Lectures on Abstract Algebra
Preliminary Version

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Part 1

Preliminaries
1. Introduction

In this section, we introduce some of the notations and definitions that we shall use throughout this book. We shall do it by investigating a few mathematical statements. Consider the following:

Statements 1.1.

1. The integer $2^{43112609} - 1$, which has 12,978,189 digits, is a prime.
2. There exist infinitely many (rational) primes.
3. Let $\pi(x) :=$ the number of positive primes $\leq x$.

Then

$$\lim_{x \to \infty} \frac{\pi(x)}{x / \log x} = 1.$$ 

4. $\sqrt{2} \notin \mathbb{Q}$, where $\mathbb{Q}$ is the set of rational numbers.
5. The real number $\pi$ is not “algebraic” over $\mathbb{Q}$.
6. There exist infinitely many real numbers (even “many more” than the number of elements in $\mathbb{Q}$) not algebraic over $\mathbb{Q}$.

To look at these statements, we need some definitions.

Let $a,b \in \mathbb{Z}$, where $\mathbb{Z}$ is the set of integers. We say that $a$ divides $b$ (in $\mathbb{Z}$) if

there exists an integer $n$ such that $b = an$, i.e., if $a \neq 0$, then $\frac{b}{a} \in \mathbb{Z}$.

We write $a \mid b$ (in $\mathbb{Z}$).

For example, $3 \mid 12$ as $12 = 4 \cdot 3$ but $5 \nmid 12$ in $\mathbb{Z}$, where $\nmid$ means does not divide.

An integer $p$ is called a prime if $p \neq 0, \pm 1$ and

$n \mid p$ in $\mathbb{Z}$ implies that $n \in \{1, -1, p, -p\}$, i.e., $n = \pm 1, \pm p$.

For example, 2, 3, 5, 7, 11, ... are prime. [The prime 2 is actually the “oddest prime of all”!] We shall need to know below that every integer $n > 1$ is divisible by some prime (to be proven later).

A complex number $x$ is called irrational if $x$ is not a rational number, i.e., $x \in \mathbb{C} \setminus \mathbb{Q} := \{ z \in \mathbb{C} \mid z \notin \mathbb{Q} \}$, where $\mathbb{C}$ is the set of complex numbers. A complex number $x$ is called algebraic (over $\mathbb{Q}$) if there exists a nonzero polynomial $f(t) = a_n t^n + a_{n-1} t^{n-1} + \cdots + a_1 t + a_0$ with $a_0, \ldots, a_n \in \mathbb{Q}$ (some $n$) not all zero satisfying $f(x) = a_n x^n +$
\[ a_{n-1}x^{n-1} + \cdots + a_1x + a_0 = 0. \] We shall write \( f \) for \( f(t) \) where \( t \) always represents a variable and \( f(x) \) means plug \( x \) into \( f \). We let

\[ \mathbb{Q}[t] := \text{the set of polynomials with rational coefficients.} \]

So \( x \) is algebraic over \( \mathbb{Q} \) if there exists \( 0 \neq f \in \mathbb{Q}[t] \) such that \( f(x) = 0 \). A complex number \( x \) that is not algebraic (over \( \mathbb{Q} \)) is called transcendental (over \( \mathbb{Q} \)), so \( x \in \mathbb{C} \) is transcendental if there exists no nonzero polynomial \( f \in \mathbb{Q}[t] \) satisfying \( f(x) = 0 \).

With the above definitions, we can look at our statements to see which ones are interesting, deep, etc.

We start with Statement 1 above. The integer 12,978,189 was shown to be prime by a UCLA team using a primality test of Lucas on Mersenne numbers. This was the first known prime to have at least ten million digits. It is, on the face of it, not very interesting. After all, it is analogous to saying 97 is a prime. However, it is interesting historically. We call an integer \( M_n := 2^n - 1 \), \( n \) a positive integer, a Mersenne number and a Mersenne prime if it is a prime. In 1644, Mersenne conjectured that for \( n \leq 257 \), the Mersenne number \( M_n \) is prime if and only if \( n \) is a prime as \( 2^{ab} - 1 \) factors if \( a, b > 1 \). It was shown in the 1880s that \( M_{61} \) is a prime. In 1903 Frank Cole showed that \( M_{67} = 193707721 \cdot 761838257287 \). [He silently multiplied out these two numbers on a blackboard at a meeting of the American Mathematical Society.] There were three more errors to Mersenne’s conjecture: \( M_{89} \) and \( M_{107} \) are prime and \( M_{257} \) is composite, i.e., not 0, \( \pm 1 \), or a prime. The correct solution of Mersenne’s conjecture is that \( M_n \) is a Mersenne prime if and only if \( n = 2, 3, 5, 7, 13, 17, 19, 31, 61, 89, 107, 127 \) was completed in 1947.

Mersenne was interested in these numbers because of their connection to a study in antiquity. An integer \( n > 1 \) is called perfect if

\[ n = \sum_{0<d<n, d|n} d \quad \text{or} \quad 2n = \sum_{0<d\leq n, d|n} d, \]

e.g., \( 6 = 1 + 2 + 3 \). The first equation says that \( n \) is the sum of its aliquot divisors. A theorem in elementary number theory says:

**Theorem 1.2.** (Euclid/Euler) An even number \( N \) is perfect if and only if \( N = \frac{1}{2}p(p + 1) \) with \( p \) a Mersenne prime (so \( N = 2^{n-1}(2^n - 1) \) with \( p = 2^n - 1 \)).

It is still an open problem whether there exist infinitely many even perfect numbers (or equivalently, infinitely many Mersenne primes). It is also an open problem whether there exist any odd perfect numbers.
1. INTRODUCTION

Statement 2 is very interesting and is due to Euclid. It may be the first deep mathematical fact that one learns. The proof is quite simple. If the result is false, let \( p_1, \ldots, p_n \) be a complete list of (positive) primes and set \( N = p_1 \cdots p_n + 1 \). We use (cf. Exercise 1.12(1)) that says:

If \( x \mid y \) and \( x \mid z \) in \( \mathbb{Z} \), then for all \( a, b \in \mathbb{Z} \), we have \( x \mid ay + bz \),

i.e., if an integer divides two integers, then it divides any \( \mathbb{Z} \)-linear combination of those two numbers. From this, we conclude that if \( p_i \mid N \) then \( p_i \mid N - p_1 \cdots p_n \), so \( p_i \mid 1 \), which is impossible. This means that \( N \) is a prime different from any of \( p_1, \ldots, p_n \), a contradiction.

Statement 3 is a deep theorem called the \textit{Prime Number Theorem} or \textit{PNT}. It quantifies Statement 2. It was first conjectured by Gauss and proven independently in 1896 by de la Vallée Poussin and Hadamard using complex analysis and fundamental work of Riemann. One shows that the \textit{zeta function}, \( \zeta(s) = \sum_{n=1}^{\infty} 1/n^s \), where \( s \) is a complex variable, has no root on the line \( \text{Re}(s) = 1 \), i.e., \( \zeta(1 + \sqrt{-1}t) \) is never zero, and that this is equivalent to the PNT. [The most famous problem in mathematics is the Riemann Hypothesis, which says the analytic continuation of the zeta function for \( \text{Re}(s) > 0 \) has all its roots on the line \( \text{Re}(s) = 1/2 \).] An “elementary” proof in number theory means that it does not use complex analysis. An elementary proof may be very difficult. An elementary proof of PNT was discovered in 1949 by Selberg and Erdös. It is much harder than the non-elementary proof.

Statement 4 is not very interesting, although the Pythagoreans knew it but were afraid to leak the knowledge, fearing catastrophe if it got out. It did. They were apparently right. To prove this, we need:

\textbf{Theorem 1.3.} (The Fundamental Theorem of Arithmetic) Every integer \( n > 1 \) is a product of positive primes unique up to order, i.e., there exist unique primes \( 1 < p_1 < \cdots < p_r \) and unique integers \( e_1, \ldots, e_r > 0 \) such that

\[ (*) \quad n = p_1^{e_1} \cdots p_r^{e_r}. \]

[We call (*) the \textit{standard representation} or \textit{standard factorization} of \( n \). We show Statement 4 assuming that this theorem is true. [Euclid knew its proof and we shall prove it in Theorem 4.16.]]

If \( \sqrt{2} \) is rational, then \( \sqrt{2} = m/n \) for some integers \( m, n \) with \( n \neq 0 \). This means that we have an equation of integers \( 2n^2 = m^2 \). But the prime two occurs on the left hand side to an odd power and on the right hand side to an even power, contradicting the uniqueness part of
the Fundamental Theorem of Arithmetic. Thus $\sqrt{2}$ is irrational, i.e., a complex number that is not rational.

Statement 5 was proven by Lindemann in 1882. This is historically interesting as it solves the famous Greek construction problem, “squaring the circle”, which asks whether one can construct a circle with the same area as a given square using only a straight-edge and compass (according to specific rules). Its proof is not easy. One must reduce this geometric problem to an algebraic one in field theory, which is then proven using analysis. A proof is given in Section 71. It is easier to show that the real number $e$ is transcendental. (Cf. Section 69.) The first real number shown to be transcendental was the number $\sum_{n=1}^{\infty} \frac{1}{10^n}$ shown to be transcendental by Liouville [cf. Section 68]. [That it is transcendental is essentially due to the fact that this series converges much too rapidly.] An even deeper result, solved independently by Gelfond and Schneider [which we do not prove], is

**Theorem 1.4.** (Hilbert’s Seventh Problem) Let $\alpha$ and $\beta$ be algebraic over $\mathbb{Q}$ with $\alpha \neq 0, 1$ and $\beta$ irrational. Then $\alpha^\beta$ is transcendental. In particular $\sqrt{2}^{\sqrt{2}}$ is transcendental.

Lindemann actually proved that if $\alpha$ is algebraic and nonzero, then $e^\alpha$ is not algebraic. Since Euler showed that $-1 = e^{\pi i}$, it follows that $\sqrt{-1}\pi$ is not algebraic; and from this it follows that $\pi$ is transcendental, using the fact that the product of algebraic numbers is algebraic, which we prove in Theorem 47.19.

Statement 6 is not as interesting as Statement 5 for the word “transcendental” is a misnomer. However, the proof of this by Cantor changed mathematics and caused much philosophical distress. Indeed, Cantor showed that there were “many more” reals than rationals. What does this mean?

Let $A$ and $B$ be two sets. [What is a set?] We say that $A$ and $B$ have the same cardinality and write $|A| = |B|$ if there exists a one-to-one and onto function, i.e., a bijection, $f : A \to B$. We call such an $f$ a bijective function. Recall that a function or map $f : A \to B$ is one-to-one or injective if

$$f(a_1) = f(a_2) \implies a_1 = a_2$$

(where $\implies$ means implies). Let $f(A) := \{f(a) \mid a \in A\}$ denote the image of $A$. Then $f$ is one-to-one is if and only if

$$f^{-1} : f(A) \to A \text{ given by } f(a) \mapsto a$$
is a function. We say $f$ is onto or surjective if $B = f(A)$. [We also write the function as $A \xrightarrow{f} B$.]

**Review 1.5.** You should review the definitions of functions, inverses, etc. that we assume you have learned.

**Examples 1.6.**
1. \{1, 2, \ldots, 26\} and \{a, b, \ldots, z\} have the same cardinality.
2. \(\mathbb{Z}\) and \(2\mathbb{Z} := \{2n \mid n \in \mathbb{Z}\}\) have the same cardinality, as $f : \mathbb{Z} \to 2\mathbb{Z}$ given by $n \mapsto 2n$ is a bijection.
3. Let \(\mathbb{Z}^+\) denote the set of positive integers and \(\mathbb{N} = \mathbb{Z}^+ \cup \{0\}\) the set of non-negative integers. Then \(\mathbb{Z}, \mathbb{N}, \) and \(\mathbb{Z}^+\) all have the same cardinality. Indeed, the maps
   $$ f : \mathbb{Z} \to \mathbb{N} \text{ by } n \mapsto \begin{cases} 2n - 1 & \text{if } n > 0 \\ -2n & \text{if } n \leq 0 \end{cases} $$

   and $g : \mathbb{N} \to \mathbb{Z}^+$ by $n \mapsto n + 1$ are both bijections.

   If \(A\) is a set, we call the symbol $|A|$ the **cardinality** of \(A\) and interpret it to mean the “number of elements” in \(A\). For example, if \(A\) is finite, i.e., the set \(A\) has finitely many elements [we write $|A| < \infty$], then $|A|$ is the number of elements in \(A\), e.g., we have $|\{1, \ldots, 26\}| = 26$. If \(A\) is not finite, we of course say that \(A\) is infinite. If a set \(A\) satisfies $|A| = |\mathbb{Z}|$, i.e., there exists a bijection $f : \mathbb{Z} \to A$ (equivalently, there exists a bijection $g : A \to \mathbb{Z}$), we say \(A\) is **countable**. (Some authors call such a set **countably infinite** or **denumerable**.)

The following facts can be shown [and we leave their proofs as exercises]:

**Facts 1.7.**
1. A subset of a countable set is either finite or countable.
2. \(\mathbb{Q}\) is countable.
3. Let \(\mathbb{C}\) be the set of complex numbers. Then \(\{\alpha \in \mathbb{C} \mid \alpha \text{ algebraic }\}\), the set of complex algebraic numbers, is countable.

**Theorem 1.8.** (Cantor) The set \(\mathbb{R}\) of real numbers is not countable.

**Proof.** Suppose that \(\mathbb{R}\) is countable. As \(\mathbb{R}\) and the closed interval \([0, 1]\) have the same cardinality by Exercise 1.12(10), the interval \([0, 1]\) must also be countable. As \(\mathbb{Z}^+ \subset \mathbb{Z}\) is not finite, it is also countable by the first fact (or by Example 1.6(3)). Therefore, there exists a bijection $f : \mathbb{Z}^+ \to [0, 1]$, i.e., the closed interval \([0, 1]\) is covered by a sequence.
Write each $f(n)$ in (base ten) decimal form, say
\[ f(1) = .a_1a_2a_3 \cdots a_n \cdots \]
\[ f(2) = .a_2a_2a_3 \cdots a_2n \cdots \]
\[ \vdots \]
\[ f(n) = .a_na_2a_3 \cdots a_nn \cdots \]
\[ \vdots \]
For each $n$ choose $b_n \in \{1, \ldots, 8\}$ with $b_n \neq a_{nn}$. [We omit 0 and 9 to prevent round off problems.] Let $b = b_1b_2b_n \cdots$. Then $b$ and $f(n)$ differ in the $n$th digit as $b_n \neq a_{nn}$. Moreover, as $b_n \neq 0,9$, the numbers $b$ and $f(n)$ cannot be the same real number. Thus $b \neq f(n)$ for all $n \in \mathbb{Z}^+$, so $f$ is not onto, a contradiction. (We have really used the fact that $\mathbb{R}$ is complete, i.e., that $b$ is a real number.) \[ \square \]

**Warning 1.9.** Sets can be tricky – again we ask what is a set?

Indeed, consider the statement:

**Statement 1.10.** There is a universe, i.e., a set that contains all other sets.

Now you should consider the following:

**Contemplate 1.11.** (Russell’s Paradox) Let
\[ A := \{ B \text{ is a set} \mid B \notin B \}. \]

Is $A \in A$? Is $A \notin A$?

What can you conclude?

**Exercises 1.12.**

1. Show if a nonzero integer $a$ divides integers $b$ and $c$, it divides $bx + cy$ for any integers $x$ and $y$.
2. Show if $a^n - 1$ is a prime with $n > 1$, then $a = 2$ and $n$ is a prime. [The converse is false as $23 \mid M_{11}$.]
3. If $2^n + 1$ is a prime, what can you say about $n$? Prove your assertion.
4. Let $a, b, n \in \mathbb{Z}$ with $n > 1$ and $b \neq 0$. Determine when $n^{\frac{b}{n}}$ is a rational number. Prove your determination. (You can use the Fundamental Theorem of Arithmetic.)
5. Assume that the real number $\sqrt{\pi}$ is transcendental. Show that $\pi$ is transcendental.
6. Let \( f : A \to B \) be a map of sets. Prove that \( f \) is injective if and only if given any set \( C \) and any two set maps \( g_i : C \to A, \ i = 1, 2, \) with compositions \( f \circ g_1 = f \circ g_2, \) then \( g_1 = g_2. \)

7. Let \( f : A \to B \) be a map of sets. Prove that \( f \) is surjective if and only if given any set \( C \) and any two set maps \( h_i : B \to C, \ i = 1, 2, \) with compositions \( h_1 \circ f = h_2 \circ f, \) then \( h_1 = h_2. \)

8. Show a subset of a countable set is either countable or finite.

9. If the sets \( A \) and \( B \) are countable, show that the cartesian product

\[
A \times B := \{(a, b) \mid a \in A, \ b \in B\}
\]

is also countable. Using induction (to be discussed), show if

\( A_1, \ldots, A_n \) are countable so is \( A_1 \times \cdots \times A_n. \)

[Hint: Show \( \mathbb{Z} \times \mathbb{Z} \) (or \( \mathbb{Z}^+ \times \mathbb{Z}^+ \) where \( \mathbb{Z}^+ := \{n \in \mathbb{Z} \mid n > 0\} \)) is countable by drawing a big (= infinite) matrix. In a similar way one proves Fact 1.7(2). In fact, the cartesian product of countable sets is countable, which is needed to prove Fact 1.7(3). Can you prove this?]

10. The closed interval \([0, 1]\) and the set of all real numbers \( \mathbb{R} \) have the same cardinality.

11. Any two (finite) line segments (so each has more than one point) have the same cardinality. [Hint: Draw the two line segments parallel to each other.]
CHAPTER I

The Integers

In this chapter, we investigate the basic properties of the set of integers. In particular, we show that every integer can be factored into a product of primes, unique up to order. To do so we introduce concepts that shall be generalized. We assume that the reader has seen proofs by induction. We begin the chapter with an equivalent form of induction called the Well-Ordering Principle and use it to establish facts about division of integers. In particular, we prove the division algorithm. Of great historical importance is the notion of prime and the property discovered by Euclid that characterizes a prime that became a cornerstone of modern algebra.

2. Well-Ordering and Induction

We denote the empty set by \( \emptyset \). In this section, we are interested in induction which you should have seen and turns out to be equivalent to the following:

The Well-Ordering Principle 2.1. Let \( \emptyset \neq S \subset \mathbb{Z}^+ \). Then \( S \) contains a least (also called a minimal) element, i.e.,

There exists an \( a \in S \) such that if \( x \in S \), then \( a \leq x \).

This is an axiom that we accept as being true. You can visualize it by drawing the positive real line and ticking off the integers. You move to the right until you get to the first element in \( S \).

One can modify the Well-Ordering Principle to the seemingly more general:

The Modified Well-Ordering Principle 2.2. Let \( \emptyset \neq T \subset \mathbb{Z} \). Suppose that there exists an \( N \in \mathbb{Z} \) such that \( N \leq x \) for all \( x \in T \), i.e., \( T \) is bounded from below. Then \( T \) contains a least element.

Proof. Let \( f : \mathbb{Z} \to \mathbb{Z} \) be the map defined by \( x \mapsto x + |N| + 1 \) where here \( |N| \) means the absolute value of \( N \). Then \( f \) is a bijection. [What is \( f^{-1} \)?] Hence the restriction of \( f \) to the subset \( T \), which we denote by \( f|_T : T \to \mathbb{Z} \), is also injective. As \( f(T) = f|_T(T) \subset \mathbb{Z}^+ \), there exists a least element \( f(a) \in f(T) \), \( a \in T \), by the Well-Ordering
Principle. Then \( a \) is the least element of \( T \) as \( f \) preserves \( \leq \). (You should check this.) □

**Remark 2.3.** In an analogous way, if \( \emptyset \neq T \subset \mathbb{Z} \) is bounded from above, i.e., there exists an integer \( N \) satisfying \( s \leq N \) for all \( s \in T \), then \( T \) contains a largest (also called a maximal) element, i.e., an element \( a \in T \) such that \( x \leq a \) for all \( x \) in \( T \).

We have the following simple applications of well-ordering.

**Application 2.4.** There exists no integer \( N \) satisfying \( 0 < N < 1 \) (or strictly between any integers \( n \) and \( n + 1 \)).

**Proof.** Let \( S = \{ n \in \mathbb{Z} \mid 0 < n < 1 \} \). If \( \emptyset \neq S \) then there exists a least element \( N \in S \) by well-ordering, so \( 0 < N < 1 \) which implies that \( 0 < N^2 < N < 1 \) – can you prove this? – and \( N^2 \in \mathbb{Z} \) contradicting the minimality of \( N \). □

**Application 2.5.** Let \( P(n) \) be a statement that is true or false depending on \( n \in \mathbb{Z}^+ \). If \( P(n) \) is not true for some \( n \) then by the Well-Ordering Principle, there exists a least element \( n \in \mathbb{Z}^+ \) such that \( P(n) \) is false. We call \( P(n) \) a minimal counterexample.

**Application 2.6.** Suppose that \( S \subset \mathbb{Z}^+ \) and \( 1 \in S \). If \( S \) satisfies the condition that whenever \( n \in S \) also \( n + 1 \in S \), then \( S = \mathbb{Z}^+ \).

**Proof.** Let \( T = \mathbb{Z}^+ \setminus S = \{ n \in \mathbb{Z}^+ \mid n \not\in S \} \), i.e., \( \mathbb{Z}^+ = T \cup S \), the union of \( T \) and \( S \), and \( \emptyset = T \cap S \), the intersection of \( T \) and \( S \). (We say that \( \mathbb{Z}^+ \) is the disjoint union of \( S \) and \( T \), and we write \( \mathbb{Z}^+ = S \lor T \).)

Suppose that \( T \neq \emptyset \). Then by well-ordering there exists a least positive element \( n \in T \). In particular, \( n - 1 \not\in T \). As \( 1 \not\in T \), \( n > 1 \), so we have \( n - 1 \in \mathbb{Z}^+ \). Hence by minimality, \( n - 1 \in S \). The hypothesis now implies that \( n \in S \), a contradiction. Thus \( T = \emptyset \) so \( \mathbb{Z}^+ = S \). □

The applications yield the following:

**Application 2.7.** (First Principle of Finite Induction) For each positive integer \( n \), let \( P(n) \) be a statement that is true or false. Suppose we know that

1. (Base Case) \( P(1) \) is true.
2. (Induction Step) If \( P(n) \) is true then \( P(n + 1) \) is true.

Then \( P(n) \) is true for all positive integers \( n \).

The hypothesis in the induction step is called the induction hypothesis.
Often the following equivalent variant of the First Principle of Finite Induction is more useful. That it is equivalent is left as an exercise.

**Application 2.8.** (Second Principle of Finite Induction) For each positive integer \( n \), let \( P(n) \) be a statement that is true or false. Suppose that the assumption that \( P(m) \) is true for all positive integers \( m < n \) implies that \( P(n) \) is true. Then \( P(n) \) is true for all positive integers \( n \).

**Remark 2.9.** If \( n = 1 \) then \( \{m < 1 \mid m \in \mathbb{Z}^+\} \) is empty, so you still must show that \( P(1) \) is true if you wish to use the Second Principle of Finite Induction, as your proof would fail for \( n = 1 \) otherwise. Remember that well-ordering needs a non-empty set.

**Remark 2.10.** Note that using the Modified Well-Ordering Principle, you can start induction at any integer and prove things from that point on, e.g., you can start at 0.

We assume that you have seen induction proofs in the past, e.g., in your linear algebra course. Induction proofs can be very complicated. We give such a complicated induction proof. Note in the proof the formal steps versus the mathematical ones.

**Example 2.11.** The product of any \( n \geq 1 \) consecutive positive integers is divisible by \( n! \).

**Proof.** Let \( m \) and \( n \) be two positive integers and set

\[
(m)_n := m(m + 1) \cdots (m + n - 1),
\]

the product of \( n \) consecutive integers starting from \( m \).

**Claim 2.12.** We have \( n! \mid (m)_n \) for all positive integers \( m \) and \( n \).

Of course, the claim is exactly what we want to prove. Note that there are two integers in the claim. We have

If \( m \in \mathbb{Z}^+ \) is arbitrary and \( n = 1 \) then \( (m)_n = m \) and \( n! = 1! \mid m \).

We assume

**Induction Hypothesis I.** The claim holds for fixed \( n = N - 1 \) and for all \( m \).

This is the induction step on \( n \), using \( N \) for clarity. As we know that this holds for \( n = 1 \) and for all \( m \), we must show that the result holds for \( N \) and all \( m \), i.e., we must show

\[
(N - 1)! \mid (m)_{N-1} \text{ for all } m \implies N! \mid (m)_N \text{ for all } m.
\]

Let \( m = 1 \). Then \( (1)_N = N! \) and \( N! \mid (1)_N \). Note that this is the first step of an induction on \( m \). So we can assume
**Induction Hypothesis II.** The claim holds for fixed $n = N$ and $m = M$. (Notationally it is easier to use $M$ rather than $M - 1$).

We must show that Induction Hypotheses I and II imply the result for $n = N$ and $m = M + 1$, i.e.,

if $N! \mid (M)_N$ and $(N - 1)! \mid (M + 1)_{N-1}$ then $N! \mid (M + 1)_N$.

Note that if we show this, then we have completed the induction step for the Induction Hypothesis II hence the result would be true for $n = N$ and for all $m$. This in turn completes the induction step for Induction Hypothesis I and hence would prove the claim.

Note also, so far even though a lot has been written, everything has been completely formal and no real work has been done except for the facts that $1 \mid m$ and $n! \mid n!$. We finally must do the mathematics. This comes from the equation

$$(M + 1)_N - (M)_N = (M + 1)(M + 2) \cdots (M + N) - M(M + 1) \cdots (M + N - 1)$$

$$= (M + 1) \cdots (M + N - 1)[(M + N) - M] \quad \text{(factoring)}$$

$$= N(M + 1) \cdots (M + N - 1) = N(M + 1)_{N-1}.$$

By Induction Hypothesis I, we have

$$(N - 1)! \mid (M + 1)_{N-1} \text{ so } N! \mid N(M + 1)_{N-1},$$

and by Induction Hypothesis II, we have $N! \mid (M)_N$. This means that

$N! \mid N(M + 1)_{N-1} + (M)_N, \text{ i.e., } N! \mid (M + 1)_N$

by Exercise 1.12(1), completing the induction step for Induction Hypothesis II, which completes the induction step for Induction Hypothesis I as needed. \(\square\)

**Corollary 2.13.** Let $n \in \mathbb{Z}^+$. Then there exist infinitely many $n$ consecutive composite (i.e., non-prime and not 0 or ±1) positive integers.

**Proof.** Let $m \in \mathbb{Z}^+$ and $N = (m)_{n+1} = m(m + 1) \cdots (m + n)$. Then by the example, we have $(n + 1)! \mid N$. If follows that for any integer $s$ satisfying $2 \leq s \leq n + 1$, we have $s \mid N + s$. As $n + 1 < N + s$, the positive integers

$$(*) \quad N + 2, N + 3, \ldots, N + n + 1$$

are all composite. \(\square\)
Of course, the same proof works with \( m = 1 \), so we did not need the example to prove this. However, (*) says there are arbitrarily large gaps between primes. On the other hand, primes are not ‘sparse’, as Chebyshev proved Bertrand’s Hypothesis, which says

If \( n \in \mathbb{Z}^+ \) then there exists a prime \( p \) satisfying \( n \leq p \leq 2n \).

Of course, \( 2n \) in the above can only be prime if \( n = 1 \). In fact, earlier Euler showed that the infinite sum \( \sum_{p \text{ a positive prime}} \frac{1}{p} \) diverges, so there are ‘many more’ primes than say squares (in some sense).

**Corollary 2.14.** Let \( m \) and \( n \) be positive integers with \( n \leq m \). Define the binomial coefficients

\[
\binom{m}{n} := \frac{m!}{(m-n)!n!} = \frac{m(m-1)\cdots(m-n+1)}{n!}
\]

and

\[
\binom{-m}{n} := (-1)^n\frac{m(m+1)\cdots(m+n-1)}{n!}.
\]

Then \( \binom{m}{n} \) and \( \binom{-m}{n} \) are integers.

**Proof.** We know that \( \binom{m}{n} = \frac{(m-n+1)n}{n!} \) is an integer by the example. The second statement follows from the identity

\[
\binom{-m}{n} = (-1)^n\binom{m+n-1}{n}.
\]

__Corollary 2.15.** Let \( p > 1 \) be a prime. Then

\[
\binom{p}{1}, \binom{p}{2}, \ldots, \binom{p}{p-1}
\]

are all divisible by \( p \), i.e., \( p \mid \binom{p}{n} \) for \( 1 \leq n \leq p-1 \).

**Proof.** If \( 1 \leq n \leq p - 1 \), then the example says that \( n! \mid p(p-1)\cdots(p-n+1) = (p-n+1)_n \).

We also know that \( s \) and \( p \) have no common non-trivial factors if \( 1 < s < p \). \([\text{Proof?]})\] We say that \( s \) and \( p \) are relatively prime. (Cf. 4.4.) We shall show below (cf. Theorem 4.10) but assume for now the following:

**Theorem 2.16.** (General Form of Euclid’s Lemma) Let \( a, b, c \) be nonzero integers with \( a \) and \( b \) relatively prime integers. If \( a \mid bc \) then \( a \mid c \).
Applying Euclid’s Lemma, we conclude that 
\[ n! \mid (p - 1) \cdots (p - n + 1) \text{ so } pn! \mid p(p - 1) \cdots (p - n + 1) \]
for all \( n \) satisfying \( 1 \leq n \leq p - 1 \). (Why?) □

**Exercises 2.17.**

1. Prove that the number of subsets of a set with \( n \) elements is \( 2^n \).
2. When Gauss was ten years old he almost instantly recognized that 
   \[ 1 + 2 + \cdots + n = \frac{n(n+1)}{2}. \] 
   [Actually, what he did was a bit harder.]
   What is a formula for the sum of the first \( n \) cubes? Prove your result.
3. The first nine Fibonacci numbers are 1, 1, 2, 3, 5, 8, 13, 21, 34. What is the \( n \)th Fibonacci number \( F_n \)? Show that \( F_n < 2^n \).
4. Note that Euclid’s proof of the infinitude of primes clearly shows that if \( p_n \) is the \( n \)th prime, then \( p_{n+1} \leq p_n^2 + 1 \). Be more careful and show that \( p_{n+1} \leq 2^{2n+3} \)? Using this, show that \( \pi(x) \geq \log \log(x) \), where \( \pi(x) \) is the number of primes less than \( x \) if \( x \geq 2 \). [This is a bad estimate.]
5. Prove that the Well-Ordering Principle, the First Principle of Finite Induction, and the Second Principle of Finite Induction are all equivalent.
6. State and prove the binomial theorem. What algebraic properties do you need for your proof to work?
7. Fill in the missing details in the proof of Corollary 2.15.

**3. Addendum: The Greatest Integer Function**

Our (double) induction showing that binomial coefficients are integers was complicated to write, although the mathematics was simple. In this addendum, we shall give another proof, mathematically more complicated, but of greater interest. We shall use some results that we shall prove in the next section, in particular, the Division Algorithm 4.2 and the Fundamental Theorem of Arithmetic 4.16. This alternate proof uses the following function:

**Definition 3.1.** The *greatest integer function* 
\[ [x] : \mathbb{R} \to \mathbb{Z} \] is given by \([x] := \text{the greatest integer } n \text{ satisfying } n \leq x\), i.e., if \( x \in \mathbb{R} \), using Application 2.4, we can write 
\[ x = n + x_0 \text{ with } 0 \leq x_0 < 1, n \in \mathbb{Z} \text{ then } [x] = n. \]

This function satisfies:
Properties 3.2. Let \( m \) be a positive integer. Then the following are true:

1. \( |x| \leq x < |x| + 1. \)
2. \( [x + m] = [x] + m. \)
3. \( \left[ \frac{x}{m} \right] = \left[ \frac{|x|}{m} \right]. \)
4. \( [x] + [y] \leq [x + y] \leq [x] + [y] + 1. \)
5. If \( n, a \in \mathbb{Z}^+ \), then \( \left[ \frac{n}{a} \right] \) is the number of integers among \( 1, \ldots, n \) that are divisible by \( a \).

Proof. (1) and (2) are easy to see.

(3): Write \( x = n + x_0 \) with \( 0 \leq x_0 < 1 \) and \( n \in \mathbb{Z} \). Using the Division Algorithm 4.2, write \( n = qm + r \) with \( q, r \in \mathbb{Z} \) and \( 0 \leq r < m \). Then we have
\[
\left[ \frac{x}{m} \right] = \left[ \frac{n + x_0}{m} \right] = \left[ \frac{qm + r + x_0}{m} \right] = q + \left[ \frac{r + x_0}{m} \right].
\]
As \( r \in \mathbb{Z} \) and \( 0 \leq r < m \), we have \( 0 \leq r \leq m - 1 \) (cf. Application 2.4); and as \( 0 \leq x_0 < 1 \), we have \( r + x_0 < m \), so
\[
\left[ \frac{x}{m} \right] = q = \left[ q + \frac{r}{m} \right] = \left[ \frac{n}{m} \right] = \left[ \frac{|x|}{m} \right].
\]

(4): Write
\[
x = n + x_0 \quad \text{with} \quad 0 \leq x_0 < 1 \quad \text{and} \quad n \in \mathbb{Z}
\]
\[
y = q + y_0 \quad \text{with} \quad 0 \leq y_0 < 1 \quad \text{and} \quad q \in \mathbb{Z}.
\]
Then \( 0 \leq x_0 + y_0 \), so
\[
[x+y] = [n+q+x_0+y_0] = n+q+[x_0+y_0] = [x]+[y]+[x_0+y_0] < [x]+[y]+2
\]
and the result follows easily.

(5): Let \( a, 2a, \ldots, ja \) denote all the positive integers \( \leq n \) and divisible by \( a \). We must show \( \left[ \frac{n}{a} \right] = j \). Clearly, \( ja \leq n < (j + 1)a \). Therefore, \( j \leq \frac{n}{a} < j + 1 \). The result now follows. \( \square \)

Let \( n \in \mathbb{Z}^+ \) and \( p > 1 \) be a prime. By the Fundamental Theorem of Arithmetic 4.16, we know that there exists a unique integer \( e \geq 0 \) such that \( p^e \mid n \) but \( p^{e+1} \nmid n \). We shall write this as \( p^e \mid\mid n \).

Theorem 3.3. Let \( n \) be a positive integer and \( p \) a positive prime. Suppose that \( p^e \mid\mid n \). Then
\[
e = \sum_{i=1}^{\infty} \left[ \frac{n}{p^i} \right].
\]
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Proof. If \( p^i > n \), then \( \left\lfloor \frac{n}{p^i} \right\rfloor = 0 \). So the sum is really a finite sum. We prove the result by induction on \( n \).

\( n = 1 \). There is nothing to prove.

We can, therefore, make the following:

**Induction Hypothesis.** Let \( e' = \sum_{i=1}^{\infty} \left\lfloor \frac{n-1}{p^i} \right\rfloor \), then \( p^e' \parallel (n-1)! \).

Let \( p^f \parallel n \). As \( n! = n(n-1)! \), by the induction hypothesis, we have \( p^e = p^e' p^f = p^{e+f} \parallel n! \), hence \( e = f + e' \) or \( f = e - e' \). So it suffices to prove that

\[
f = \sum_{i=1}^{\infty} \left\lfloor \frac{n}{p^i} \right\rfloor - \sum_{i=1}^{\infty} \left\lfloor \frac{n-1}{p^i} \right\rfloor.
\]

**Claim 3.4.** \( \left\lfloor \frac{n}{p^i} \right\rfloor - \left\lfloor \frac{n-1}{p^i} \right\rfloor = \begin{cases} 1 & \text{if } p^i \parallel n \\ 0 & \text{if } p^i \not\parallel n. \end{cases} \)

Suppose that \( p^i \parallel n \), then \( n = p^i u \) with \( u \in \mathbb{Z} \). As \( 0 < \frac{1}{p^i} < 1 \), we have

\[
\left\lfloor \frac{n}{p^i} \right\rfloor = u \text{ and } \left\lfloor \frac{n-1}{p^i} \right\rfloor = \left\lfloor u - \frac{1}{p^i} \right\rfloor = u - 1
\]
as needed.

Suppose that \( p^i \not\parallel n \). The result is immediate if \( p^i > n \), so suppose not. Then we can write \( n = p^i u + r \) with \( u \) and \( r \) integers with \( r \) satisfying \( 0 \leq r < p^i \) using the Division Algorithm 4.2. Hence

\[
\frac{n}{p^i} = u + \frac{r}{p^i}, \text{ so } \left\lfloor \frac{n}{p^i} \right\rfloor = u
\]

and

\[
\frac{n-1}{p^i} = \frac{p^i u + r - 1}{p^i} = u + \frac{r - 1}{p^i} \text{ for } 0 \leq r - 1 < p^i,
\]

so we also have \( \left\lfloor \frac{n-1}{p^i} \right\rfloor = u \) as needed. This proves the claim. But the claim means that

\[
\sum_{i=1}^{\infty} \left\lfloor \frac{n}{p^i} \right\rfloor - \sum_{i=1}^{\infty} \left\lfloor \frac{n-1}{p^i} \right\rfloor = \sum_{p^i \mid n} 1 = f. \quad \Box
\]
Example 3.5. Using the properties of \([\ ]\), we compute \(e\) such that \(7^e \mid | 1000!.

\[
\begin{align*}
\left\lfloor \frac{1000}{7} \right\rfloor & = 142 \\
\left\lfloor \frac{1000}{7^2} \right\rfloor & = \left\lfloor \frac{142}{7} \right\rfloor = 20 \\
\left\lfloor \frac{1000}{7^3} \right\rfloor & = \left\lfloor \frac{20}{7} \right\rfloor = 2 \\
\left\lfloor \frac{1000}{7^4} \right\rfloor & = \left\lfloor \frac{2}{7} \right\rfloor = 0,
\end{align*}
\]

so \(e = 142 + 20 + 2 = 164\), i.e., \(7^{164} \mid 1000!\).

We can generalize our result about binomial coefficients to show that multinomial coefficients are integers.

Corollary 3.6. Suppose that \(a_1, \ldots, a_r\) are non-negative integers satisfying \(a_1 + \cdots + a_r = n\). Then the multinomial coefficient \(\frac{n!}{a_1! \cdots a_r!}\) is an integer.

Proof. By the Fundamental Theorem of Arithmetic 4.16 and the theorem, it suffices to prove that for each prime \(p\), we have

\[
\sum_{i=1}^{\infty} \left\lfloor \frac{n}{p^i} \right\rfloor \geq \sum_{i=1}^{\infty} \left\lfloor \frac{a_1}{p^i} \right\rfloor + \cdots + \sum_{i=1}^{\infty} \left\lfloor \frac{a_r}{p^i} \right\rfloor.
\]

Applying Property 3.2(4) shows that

\[
\left\lfloor \frac{a_1}{p^i} \right\rfloor + \cdots + \left\lfloor \frac{a_r}{p^i} \right\rfloor \leq \left\lfloor \frac{a_1 + \cdots + a_r}{p^i} \right\rfloor = \left\lfloor \frac{n}{p^i} \right\rfloor
\]

for all \(i\). Summing over \(i\) yields the result. \(\square\)

That the product of \(n\) consecutive positive integers is divisible by \(n!\) and that binomial coefficients are integers is now easy to see as

\[
\frac{(m)_n}{n!} = \binom{m+n-1}{n}
\]

for all \(m \in \mathbb{Z}^+\) and

\[
\binom{a}{b} = \frac{a!}{b!(a-b)!}
\]

for all \(a, b \in \mathbb{Z}^+\) with \(b \leq a\).

Chebyshev cleverly used the binomial coefficient \(\binom{2n}{n} = \frac{(2n)!}{n!n!}\). For each positive prime \(p \leq 2n\), let \(r_p \in \mathbb{Z}\) satisfy

\[
p^{r_p} \leq 2n < p^{r_p+1}
\]
and \(p^e \parallel (\binom{2n}{n})\). Then

\[
e = \sum_{i=1}^{\infty} \left\lfloor \frac{2n}{p^i} \right\rfloor - 2 \sum_{i=1}^{\infty} \left\lfloor \frac{n}{p^i} \right\rfloor.
\]

We have

\[
\left\lfloor \frac{2n}{p^i} \right\rfloor = \left\lfloor \frac{n}{p^i} \right\rfloor \leq \left\lfloor \frac{n}{p^{i-1}} \right\rfloor + 1
\]

by Property 3.2(4), so \(e \leq \sum_{i=1}^{r_p} 1 = r_p\). This is used to prove

(3.7) \(\binom{2n}{n} = \frac{(2n)!}{n!n!} \) divides \(\prod_{p \leq 2n} p^{r_p}\)

and if further \(n < p \leq 2n\), then \(p \mid (2n)!\) but \(p \not\mid (n!)^2\). Therefore, we see

(3.8) \(n^{\pi(2n) - \pi(n)} \leq \prod_{n < p \leq 2n} p \leq \binom{2n}{n} \leq \prod_{p \leq 2n} p^{r_p} \leq (2n)^{\pi(2n)}\).

Using this Chebyshev showed that there were positive real numbers \(c_1\) and \(c_2\) satisfying

\[
c_1 \frac{x}{\log x} \leq \pi(x) \leq c_2 \frac{x}{\log x}
\]

(cf. The Prime Number Theorem). He also showed if \(n \geq 3\) and \(2n/3 < p < n\), with \(p\) a prime, then \(p \not\mid \binom{2n}{n}\) and used this to prove Bertrand’s Hypothesis that there always exists a prime \(p\) such that \(n \leq p \leq 2n\).

**Exercises 3.9.**
1. Let \(a\) and \(b\) be positive integers. Show that \(\frac{(ab)!}{a!(b!)^a}\) is an integer.
2. Verify equations 3.7 and 3.8.

4. Division and the Greatest Common Divisor

We turn to further applications of well-ordering. We shall need the following:

**Properties 4.1.** For integers \(r, n,\) and \(m\) the following division properties hold:

1. If \(r \mid m\) and \(r \mid n\), then \(r \mid am + bn\) for all integers \(a\) and \(b\).
2. If \(r \mid n\), then \(r \mid mn\).
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(3) If \( r \mid n \) and \( n \neq 0 \), then the absolute value \( |n| \geq |r| \geq r \).
(4) If \( m \mid n \) and \( n \neq 0 \), then \( n = \pm m \).
(5) If \( mn = 0 \), then \( m = 0 \) or \( n = 0 \).
(6) If \( mr = nr \), then either \( m = n \) or \( r = 0 \).

**Proof.** Exercise, but do not use the Fundamental Theorem of Arithmetic. □

The most important elementary property about the integers is the next result, which intertwines multiplication and addition of integers. In the sequel, we shall investigate when it holds in more general situations. As it is one of our first proofs about integers, we shall put in a lot of detail. Notice that one can often assume properties about integers that we have not shown and perhaps are not clear. For example, in the proof below, we assume the transitivity of \( > \). Can this be proved or must it be axiomized?

**Theorem 4.2.** (Division Algorithm) Let \( m \) and \( n \) be integers with \( m \) positive. Then there exist unique integers \( q \) and \( r \) satisfying:

\begin{itemize}
  \item[(i)] \( n = qm + r \).
  \item[(ii)] \( 0 \leq r < m \)
\end{itemize}

[Of course, (ii) is the crux.]

**Proof.** We have two things to show: existence and uniqueness. We first show

**Uniqueness:** Let \((q, r)\) and \((q', r')\) be two pairs of integers satisfying the conclusion. We must show \( q = q' \) and \( r = r' \). We have

\[ \begin{align*}
qm + r &= n = q'm + r' \\
0 &\leq r, r' < m
\end{align*} \]

Without loss of generality, we may assume that \( r \leq r' \). By (*) , we have

\[ 0 \leq r' - r = (q - q')m. \]

If \( q - q' = 0 \), i.e., \( q = q' \), then \( r' - r = 0 \) and \( r' = r \) and we are done. So we may assume that \( q - q' \neq 0 \), i.e., \( q \neq q' \). Then \( r' - r = (q - q')m \neq 0 \) by Property 4.1(5) as \( m \neq 0 \). Therefore, we have

\[ r' - r > 0 \text{ and } m \mid r' - r. \]

Thus by Property 4.1(3), we have

\[ m \leq r' - r < r' < m, \]

a contradiction. So \( q - q' = 0 \) as needed.

[Note that we have used the important Property 4.1(5):

\[ ab = 0 \implies a = 0 \text{ or } b = 0, \]
which is equivalent to
\[ ab = ac \implies a = 0 \text{ or } b = c. \]

**Existence:**

**Case 1.** \( n > 0; \)

Intuitively we know that the number \( n \) must lie in some half open interval \([qm, (q + 1)m)\), i.e., there exists an integer \( q \) such that \( qm \leq n < (q + 1)m \). Let us show this rigorously. Let
\[ S = \{ s \in \mathbb{Z}^+ \mid sm > n \} \subset \mathbb{Z}^+. \]
As \( m > 0 \), we have \( m \geq 1 \). (Recall we have shown there is no integer properly between 0 and 1.) So \((n + 1)m = mn + m \geq n + m > n. [Can you show that \( mn > n \)?] Thus \( n + 1 \in S \) so \( S \neq \emptyset \). By the Well-Ordering Principle, there exists a least integer \( q + 1 \geq 1 \) which means \( qm \leq n < (q + 1)m \). Let \( r = n - qm \geq 0 \). We then have
\[ 0 \leq r = n - qm < (q + 1)m - qm = m, \]
so these \( q, r \) work.

**Case 2.** \( n < 0; \)

By Case 1, there exist integers \( q_1 \) and \( r_1 \) satisfying
\[ -n = |n| = q_1 m + r_1 \quad \text{and} \quad 0 \leq r_1 < m. \]
If \( r_1 = 0 \), then \( q = -q_1 \) and \( r = 0 \) work. If \( r_1 > 0 \), then
\[
\begin{align*}
n &= -q_1m - r_1 = -q_1m - m + m - r_1 \\
&= (-q_1 - 1)m + (m - r_1).
\end{align*}
\]
Since \( 0 \leq m - r_1 < m \), we have \( q = -q_1 - 1 \) and \( r = m - r_1 \) work. \( \square \)

**Question 4.3.** What if \( m < 0 \) in the above?

As mentioned above, later we shall be interested in generalizing the Division Algorithm for integers. You already have used an analogue of it in the case of the division of polynomials with real coefficients in the real numbers. Does this analogue hold if we only allow integer coefficients?

We turn to another important property about the integers.

**Definition 4.4.** Let \( n \) and \( m \) be integers with at least one nonzero. An integer \( d \) is called a **greatest common divisor** or **gcd** of \( m \) and \( n \) if \( d \) satisfies the following:

- \((i)\) \( d > 0 \).
- \((ii)\) \( d \mid m \) and \( d \mid n \).
- \((iii)\) If \( e \) is an integer satisfying \( e \mid m \) and \( e \mid n \) then \( e \mid d \).
If \( d = 1 \) is a gcd of \( n \) and \( m \), we say that \( n \) and \( m \) are relatively prime.

**Theorem 4.5.** Let \( n \neq 0 \) and \( m \) be integers. Then a gcd of \( m \) and \( n \) exists and is unique.

**Proof.** **Uniqueness:** If both \( d \) and \( d' \) satisfy (i), (ii), and (iii), then \( d \mid d' \) and \( d' \mid d \). Therefore, \( d = \pm d' \) by Property 4.1(4), hence \( d = |d| = |d'| = d' \) by (i).

**Existence:** Let

\[
S = \{am + bn \mid am + bn > 0 \text{ with } a, b \in \mathbb{Z}\}.
\]

[This is the tricky part of this proof. Where did this come from?]

As \( n \neq 0 \), we have \( |n| \in S \), so \( S \neq \emptyset \). By the Well-Ordering Principle, there exists a least element \( d = am + bn \) in \( S \) for some integers \( a, b \).

**Claim.** This \( d \) works, i.e., \( d \) satisfies (i), (ii), and (iii).

By choice, \( d \) satisfies (i). To show (ii), we use the Division Algorithm to produce integers \( q \) and \( r \) satisfying \( n = qd + r \) with \( 0 \leq r < d \). We show \( r = 0 \), i.e., \( d \mid n \). As

\[
n = dq + r = (am + bn)q + r, \quad \text{we have} \quad 0 \leq r = (1 - bq)n + (-aq)m.
\]

This means that either \( r \in S \) or \( r = 0 \). Since \( r < d \), the minimality of \( d \) implies that \( r \notin S \), so \( r = 0 \) and \( d \mid n \). Similarly, \( d \mid m \).

As for (iii), suppose that the integer \( e \) satisfies \( e \mid m \) and \( e \mid n \). Then by Property 4.1 (1), we have \( e \mid am + bn = d \). \( \Box \)

**Notation 4.6.** If \( m \) and \( n \) are integers with at least one nonzero, we shall denote their gcd by \((m, n)\).

Note that the proof even gives the following

**Bonus.** Let \( m \) and \( n \) be integers, at least one nonzero, and \( d = (m, n) \). Then there exist integers \( x \) and \( y \) satisfying

\[
d = (m, n) = mx + ny.
\]

In particular, the Diophantine equation

\[
(m, n) = mX + nY
\]

(with \( X \) and \( Y \) variables) has a solution, where the word Diophantine means that we only want integer solutions, e.g., the Diophantine equation \( X^2 = 2 \) has no solution.
**Warning 4.7.** In general, the \( x \) and \( y \) in the Bonus are not unique, e.g.,

\[
(2, 4) = 2 = 1 \cdot 4 - 1 \cdot 2 = 0 \cdot 4 + 1 \cdot 2.
\]

\[
(2, 3) = 1 = 1 \cdot 3 - 1 \cdot 2 = -1 \cdot 3 + 2 \cdot 2.
\]

See Exercise 4.17(7) for the general solution to the Diophantine equation \( d = mX + nY \).

Unfortunately, our proof does not show how to find any solution to the Diophantine equation \( (m, n) = mX + nY \). Euclid knew how to do so, and his proof constructs such a solution.

**Theorem 4.8.** (Euclidean Algorithm) Let \( a \) and \( b \) be positive integers with \( b \nmid a \). [If \( b \mid a \), we have \( (a, b) = b \).] Then there exists an integer \( k > 0 \) and equations

\[
\begin{align*}
a &= bq_1 + r_1 & \text{with} & \quad 0 < r_1 < b \\
b &= r_1q_2 + r_2 & \text{with} & \quad 0 < r_2 < r_1 \\
r_1 &= r_2q_3 + r_3 & \text{with} & \quad 0 < r_3 < r_2 \\
& \vdots \\
r_{k-2} &= r_{k-1}q_k + r_k & \text{with} & \quad 0 < r_k < r_{k-1} \\
r_{k-1} &= r_kq_{k+1}
\end{align*}
\]

for some integers \( q_1, \ldots, q_{k+1} \) and \( r_1, \ldots, r_k \).

**Proof.** Exercise 4.17(4). \( \square \)

Moreover, \( r_k = (a, b) \) and the Diophantine equation \( (a, b) = aX + bY \) has a solution.

Finding the gcd of two integers is much easier and faster than factoring numbers. Indeed, using the internet to make purchases or send confidential information is (relatively) safe because of the difficulty in factoring large numbers.

**Properties 4.9.** Let \( a, b, \) and \( c \) be integers with \( a \neq 0 \) and \( d = (a, b) \).

1. We have \( d = 1 \) if and only if \( 1 = ax + by \) for some integers \( x, y \).
2. \( \left( \frac{a}{d}, \frac{b}{d} \right) = 1 \). [Note that \( \frac{a}{d} \) and \( \frac{b}{d} \) are integers.]
3. If \( d = 1 \) and \( a \mid bc \) (in \( \mathbb{Z} \)), then \( a \mid c \).
4. If \( a \mid bc \) (in \( \mathbb{Z} \)), then \( \frac{a}{d} \mid c \).
5. If \( m > 0 \), then \( (ma, mb) = md \).

We have seen Property 4.1(3) above. As it is so important, we write it one more time.
Theorem 4.10. (General Form of Euclid’s Lemma) Let $a$, $b$ be relatively prime integers with $a$ nonzero. If $a \mid bc$ with $c$ an integer, then $a \mid c$.

A special case of this theorem is:

Theorem 4.11. (Euclid’s Lemma) Let $a$, $b$ be integers and $p$ be a prime satisfying $p \mid ab$. Then $p \mid a$ or $p \mid b$.

Of course this means

Consequence 4.12. If a prime $p$ satisfies $p \mid a_1 \cdots a_r$, with $a_i$ integers for $i = 1, \ldots, r$, then there exists an $i$ such that $p \mid a_i$.

Proof. (of the properties and Euclid’s Lemma.)

(1): We already know if $d = 1$ then $1 = ax + by$ for some $x, y \in \mathbb{Z}$. Conversely, suppose that $1 = ax + by$ for some $x, y \in \mathbb{Z}$. As $d \mid a$ and $d \mid b$, we have $d \mid 1$ by Property 4.1(1). Since $d > 0$, we have $d = 1$.

(2): By the Bonus, there exist integers $x$ and $y$ satisfying $d = ax + by$ hence $1 = \frac{a}{d}x + \frac{b}{d}y$ so (2) follows by (1).

(3): Whenever you can write 1 or 0 in a non-trivial way do so! It is a very useful technique. So by (1), we have

\[ 1 = ax + by \quad \text{for some } x, y \in \mathbb{Z} \quad \text{(Key Observation)} \]

and given such an equation, you can always multiply it, say by the integer $c$, to get

\[ c = c ax + cby. \quad \text{(Key Trick)} \]

As $a \mid a$ and $a \mid cb$, we conclude that $a \mid c$.

(4): The hypothesis means that $bc = an$ for some integer $n$, so $\frac{a}{d}n = \frac{b}{d}c$. By (2), we know that $\left(\frac{a}{d}, \frac{b}{d}\right) = 1$, hence $\frac{a}{d} \mid c$ by (3).

(5): We leave this is an exercise.

(proof of) Euclid’s Lemma. If $p$ is a prime and $p \mid ab$ but $p \nmid a$, then $(p, a) = 1$ as only $\pm 1, \pm p$ divide $p$. Since $p \mid ab$, we conclude that $p \mid b$ by (3). \qed

The converse of Euclid’s Lemma is also true, viz.,

Proposition 4.15. Let $p$ be an integer with $|p| > 1$. Then $p$ is a prime if and only if whenever $p \mid ab$, with $a$ and $b$ integers, then $p \mid a$ or $p \mid b$.

This proposition, which we leave as an exercise (cf. Exercise 4.17(8)) is a key to much of (commutative) ring theory. It says that we have two ways to define the notion of a prime element in the integers. For
more general structures, called rings, we can look at both of these conditions. Unfortunately, they need not be equivalent and it turns out that the more useful condition is the condition $p \mid ab$ then $p \mid a$ or $p \mid b$. In fact, its generalization has major repercussions not only in algebra but also in geometry. We shall investigate some of these later on.

We had a further loose end to establish, viz., a proof of the Fundamental Theorem of Arithmetic, which we turn to next. We first state it again.

**Theorem 4.16.** (The Fundamental Theorem of Arithmetic) Every integer $n > 1$ is a product of positive primes unique up to order, i.e., there exist unique primes $1 < p_1 < \cdots < p_r$ and integers $e_1, \ldots, e_r > 0$ such that (*)

\[ n = p_1^{e_1} \cdots p_r^{e_r}. \]

[We call (*) the standard representation or standard factorization of $n$, and this representation is unique.]

**Proof.** **Existence:** Let

\[ S = \{ n > 1 \text{ in } \mathbb{Z} \mid n \text{ is not a product of primes} \}. \]

We must show $S = \emptyset$. Suppose this is false. By the Well-Ordering Principle, there exists a minimal element $n \in S$. Clearly, no prime lies in $S$, so $n$ is not a prime. Hence there exist integers $n_1, n_2$ satisfying $n = n_1 n_2$ and $1 < n_i < n$, $i = 1, 2$ (as there exists an integer $1 < n_1 < n$ dividing $n$). By minimality, $n_1, n_2 \notin S$, so each is a product of primes. Hence so is $n = n_1 n_2$, a contradiction.

[This is a wonderful argument, and we shall see it quite often.]

**Uniqueness:** Suppose that

\[ p_1^{e_1} \cdots p_r^{e_r} = n = q_1^{f_1} \cdots q_s^{f_s} \]

with $1 < p_1 < \cdots < p_r$ and $1 < q_1 < \cdots < q_s$ primes and all the $e_i, f_j$ positive integers. We may assume that $p_1 \leq q_1$. As $p_1 \mid n = q_1^{f_1} \cdots q_s^{f_s}$, we must have $p_1 \mid q_i$ for some $i$ by Euclid’s Lemma. But $p_1 \leq q_i$ are primes, so we must have $i = 1$ and $p_1 = q_1$. Thus, by cancellation,

\[ p_1^{e_1-1} \cdots p_r^{e_r} = q_1^{f_1-1} \cdots q_s^{f_s} \]

and we are done by induction. What is the induction hypothesis and why are we done? \(\square\)

The proof of the existence statement in the Fundamental Theorem of Arithmetic is used repeatedly in mathematics and its consequence occurs much more frequently than the uniqueness statement. If we are studying a class of objects in which there is a set of basic objects (atoms), the basic question that arises is whether we can break up an
element in the class into a finite number of these atoms (in some way).
In the Fundamental Theorem of Arithmetic the atoms are the primes.

**Exercises 4.17.**

1. Prove all the properties in Properties 4.1 without using the Fundamental Theorem of Arithmetic.

2. Let $F = \mathbb{R}, \mathbb{C}$ or $\mathbb{Q}$ [or any field, or even any ring, cf. Definition 8.3]. Let $F[t]$ be the set of polynomials with coefficients in $F$ with the usual addition and multiplication. [Which are?] State and prove the analogue of the Division Algorithm for integers. [Use your knowledge of such division. Use the degree of a polynomial as a substitute for statement (ii) in the Division Algorithm for integers. (One usually does not define the degree of the zero polynomial, so if $r$ is the remainder, write $r = 0$ or ... ).] What can you do if you take polynomials with coefficients in $\mathbb{Z}$?

3. Prove the following modification of the Division Algorithm: If $m$ and $n$ are two integers with $m$ nonzero, then there exist unique integers $q$ and $r$ satisfying

$$n = mq + r \quad \text{with} \quad -\frac{1}{2}|m| < r \leq \frac{1}{2}|m|.$$ 

Moreover, if $m > 0$ is odd, then we can find unique integers $q$ and $r$ satisfying

$$n = mq + r \quad \text{with} \quad 0 \leq |r| < \frac{m}{2}.$$ 

[Note that in this case, $m/2$ is not an integer.]

4. Prove the Euclidean Algorithm 4.8, i.e., show that such a $k$ exists and that the Diophantine equation $(m,n) = mX + nY$ has a solution using the Euclidean Algorithm.

5. In the Euclidean Algorithm show each nonzero remainder $r_i$ ($i \geq 2$, where we view $b$ as $r_0$) satisfies $r_i < \frac{1}{2}r_{i-2}$. Deduce that the number of steps in the algorithm is less than

$$\frac{2 \log b}{\log 2},$$

where $b$ is the larger of the numbers $a, b$. Using the modified Division Algorithm, show the number of steps there is at most $\frac{\log b}{\log 2}$.

6. Solve the Diophantine equation

$$(39493, 19853) = 39493X + 19853Y.$$
7. Let \( m, n, \) and \( d \) be nonzero integers. Show that the Diophantine equation \( d = mX + nY \) has a solution if and only if \((m, n) | d\). If this is the case and \((x, y)\) is a solution in integers, then the general solution is \((x + \frac{n}{(m, n)}k, y - \frac{m}{(m, n)}k)\) with \(k \in \mathbb{Z}\).

8. Let \( p \) be an integer. Then \( p \) is a prime if and only if whenever \( p | ab \), with \( a \) and \( b \) integers, then \( p | a \) or \( p | b \).

9. Carefully write up the induction step and its proof for the uniqueness part of the Fundamental Theorem of Arithmetic in two different ways, i.e., by inducting on two different things.

10. Let \( n \) and \( m \) be integers with at least one nonzero. An integer \( l \) is called a **least common multiple** or **lcm** of \( m \) and \( n \) if \( l \) satisfies the following:

    (i) \( l \geq 0 \).
    
    (ii) \( m | l \) and \( n | l \).
    
    (iii) If \( k \) is an integer satisfying \( m | k \) and \( n | k \), then \( l | k \).

Show that an lcm of \( m \) and \( n \) exists and is unique, denote it by \([m, n]\). In addition, show that \( mn = [m, n](m, n)\).

11. Define \( \sigma : \mathbb{Z}^+ \to \mathbb{Z}^+ \) by \( \sigma(n) = \sum_{d|n} d \), the sum of the (positive) divisors of \( n \). Show

    (i) If \( m \) and \( n \) are relatively prime (positive) integers, then \( \sigma(mn) = \sigma(m)\sigma(n) \).
    
    (ii) If \( p \) is a (positive) prime integer and \( n \) an integer, then \( \sigma(p^n) = (p^{n+1} - 1)/(p - 1) \).

12. In the notation of the previous result, a positive integer \( n \) is called a **perfect number** if \( \sigma(n) = 2n \). Prove Theorem 1.2.
CHAPTER II

Equivalence Relations

This chapter is the most important foundational chapter in Part One. It introduces the notion of an equivalence relation and determines equivalent formulations of it. In particular, an equivalence relation is shown to be the same as a surjective map. The importance of equivalence relations is that they lead to ‘quotients’. This is probably the hardest concept to understand in Part One and Part Two. It is used to coarsen problems that may be more easily be solved, and hopefully leads to the solution of specific problems in which we are interested. The basic difficulty is that it may easy to define an equivalence relation on collection of sets, but it probably is not useful if it must reflect additional properties (e.g., algebraic properties) of those maps, i.e., if the surjective map arising from the equivalence relation (or defining it) does not reflect these additional properties. For example, what information does a surjective linear transformation of vector spaces give us?

5. Equivalence Relations

In this section, we study one of the basic concepts in mathematics, equivalence relations. You should have seen this concept. Recall the following:

**Definition 5.1.** Let $A$ and $B$ be two sets. A relation on $A$ and $B$ is a subset $R \subseteq A \times B$. It is customary to write $aRb$ if $(a, b) \in R$, and we shall always do so. If $A = B$ then we call such a relation a relation on $A$.

**Question 5.2.** What kind of relation is a function $f : A \to B$?

Besides functions, we are mostly interested in equivalence relations, which we also recall.
Definition 5.3. A relation $R$ on $A$ is called an equivalence relation on $A$ if for all $a, b, c$ in $A$,

(1) $aRa$. \hfill (Reflexivity)

(2) If $aRb$ then $bRa$. \hfill (Symmetry)

(3) If $aRb$ and $bRc$ then $aRc$. \hfill (Transitivity)

Remark 5.4. We often denote an equivalence relation by $\sim$ or $\approx$.

Question 5.5. Why do we need (1) in the definition, i.e., why don’t (2) and (3) imply (1)?

Whenever a new term is defined, you should try to find examples. It is usually difficult to have an intuitive idea of an abstract concept without examples that you know. We give a few examples of sets having equivalence relations on them.

Examples 5.6. The following are equivalence relations:

1. Any set $A$ under $=$ where if $a, b \in A$ then $a = b$ in $A$ means that we have equality of sets $\{a\} = \{b\}$.

2. Triangles in $\mathbb{R}^2$ under congruence (respectively, similarity).
   [Question: Is the mirror image of a triangle congruent to the original triangle?]

3. The set $\mathbb{Z} \times (\mathbb{Z} \setminus \{0\})$ under $\sim$ where
   
   $$(a, b) \sim (c, d) \text{ if } ad = bc \text{ (in } \mathbb{Z}).$$

4. $\mathbb{Z}$ under $\equiv \mod 2$ where

   $$m \equiv n \mod 2 \text{ if } m - n \text{ is even, i.e., } 2 \mid m - n.$$  

   Equivalently, if both $m$ and $n$ are odd or both are even in this example.

5. Let $R$ be $\mathbb{Q}$, $\mathbb{R}$, $\mathbb{C}$, or any field (or, in fact, any ring, e.g., $\mathbb{Z}$, cf. Definition 8.3). Set

   $$\mathcal{M}_n(R) := \{n \times n \text{ matrices with entries in } R\}.$$  

   Then $\sim$ is an equivalence relation on $\mathcal{M}_n(R)$, where

   $A \sim B$ if there exists $C \in \mathcal{M}_n(R)$ invertible such that $A = CBC^{-1}$.

   We call this equivalence relation similarity of matrices. We call two matrices $A$ and $B$ similar if $A \sim B$. [Cf. Change of Basis Theorem in linear algebra.]
6. The Change of Basis Theorem in linear algebra, in fact, leads to the following equivalence relation. Let \( R \) be any field (or ring). Define

\[ R^{m \times n} := \{ m \times n \text{ matrices with entries in } R \}. \]

Then \( \approx \) is an equivalence relation on \( R^{m \times n} \), where

\[ A \approx B \text{ if there exists invertible matrices } C \in M_m(R) \text{ and } D \in M_n(R) \text{ satisfying } A = CBD. \]

We say two matrices \( A \) and \( B \) in \( R^{m \times n} \) are equivalent if \( A \approx B \).

7. Let \( R \) be a field (or ring). Then \( \sim \) is an equivalence relation on \( M_n(R) \), where

\[ A \sim B \text{ if there exists } C \in M_n(R) \text{ invertible such that } A = CBC^t, \]

where \( C^t \) is the transpose of \( C \).

8. Define \( \sim_u \) on \( M_n(\mathbb{C}) \) by

\[ A \sim_u B \text{ if there exists } C \in M_n(\mathbb{C}) \text{ invertible such that } A = CBC^*, \]

where \( C^* \) is the adjoint of \( C \). Then \( \sim_u \) is an equivalence relation on \( M_n(\mathbb{C}) \).

**Definition 5.7.** Let \( \sim \) be an equivalence relation on a set \( A \). For each \( a \) in \( A \), the set

\[ [a] = [a]_{\sim} := \{ b \in A \mid a \sim b \} \]

is called the equivalence class of a relative to \( \sim \). We shall usually write \( \pi \) for \([a]\)

and call the following set of subsets of \( A \),

\[ \overline{A} = A/\sim := \{ \pi \mid a \in A \}, \]

the set of equivalence classes of \( \sim \) on \( A \).

**Warning 5.8.** The equivalence class of \( a \) is a subset of \( A \) not (in general) an element of \( A \), i.e., \( \overline{a} \subset A \), not \( \overline{a} \in A \).

Let \( \sim \) be an equivalence relation on \( A \). Then we have a map

\[ \overline{-} : A \rightarrow \overline{A} \text{ given by } a \mapsto \overline{a}. \]

This map is clearly surjective and is called the natural or canonical surjection. [It is the “obvious” map as it depends on no choices.] Thus if \( \alpha \) is an element in \( \overline{A} \), there exists an element \( a \) in \( A \) such that \( \alpha = \overline{a} \).
Examples 5.9. 1. Let $\sim$ on $\mathbb{Z} \times (\mathbb{Z} \setminus \{0\})$ be given by 

$$(a, b) \sim (c, d) \text{ if } ad = bc \in \mathbb{Z}.$$ 

Then we have 

$$\frac{(a, b)}{(a, b)} \text{ is just } \frac{a}{b} \text{ in } \mathbb{Q} = \mathbb{Z} \times (\mathbb{Z} \setminus \{0\}) / \sim.$$ 

2. Consider $\mathbb{Z}$ under $\equiv \mod 2$. Then we have 

$${\{\text{all even integers}\}} = 0 = \ldots = \pm 2n = \ldots.$$ 

$${\{\text{all odd integers}\}} = 1 = \ldots = \pm (2n + 1) = \ldots.$$ 

We shall write $\mathbb{Z} = \mathbb{Z} / \equiv \mod 2$ as $\mathbb{Z} / 2\mathbb{Z}$, so 

$$\mathbb{Z} / 2\mathbb{Z} = \{\overline{0}, \overline{1}\}.$$ 

Recall the following:

Definition 5.10. Let $A_i, i \in I$, be sets. [We call $I$ an indexing set.] The union of the sets $A_i$ is the set 

$$\bigcup_{i \in I} A_i := \{x \mid \text{there exists an } i \in I \text{ such that } x \in A_i\}.$$ 

We shall usually denote this set by $\bigcup_I A_i$. This union is a disjoint union if $A_i \cap A_j = \emptyset$ for all $i, j \in I$ with $i \neq j$. If this is the case, we denote it by $\bigvee_I A_i$. Of course, the intersection of the sets $A_i$ is the set 

$$\bigcap_{i \in I} A_i := \{x \mid x \in A_i \text{ for all } i \in I\},$$ 

and we denote by $\bigcap_I A_i$.

Proposition 5.11. Let $\sim$ be an equivalence relation on $A$. Then 

$$A = \bigvee_{\overline{a}} \overline{a}.$$ 

In particular, if $a, b \in A$, then 

either $\overline{a} = \overline{b}$ or $\overline{a} \cap \overline{b} = \emptyset$, 

hence 

$$\overline{a} = \overline{b} \text{ if and only if } a \sim b.$$ 

Proof. As $a \sim a$ for all $a \in A$ by Reflexivity and $a \in \overline{a}$ by definition, we have $A = \bigcup_a \overline{a}$. By Symmetry, we have $a \sim b$ if $\overline{a} = \overline{b}$. If $c \in \overline{a} \cap \overline{b}$, then $c \sim a$ and $c \sim b$, so $a \sim c$ by Symmetry hence $a \sim b$ by Transitivity, so $a \in \overline{b}$, i.e., $\overline{a} \subset \overline{b}$. Similarly, we have $\overline{b} \subset \overline{a}$, hence $\overline{a} = \overline{b}$. \qed
Definition 5.12. Let $\sim$ be an equivalence relation on $A$. An element $x \in \overline{a}$ is called a representative of $\overline{a}$. For example, $a$ is a representative of $\overline{a}$. So by Proposition 5.11, if $x$ is a representative of $\overline{a}$, then $x = a$. A system of representatives for $A$ relative to $\sim$ is a set

$$S = \{\text{precisely one element from each equivalence class}\},$$

so $A = \bigvee_S \overline{x}$. For example, if $\sim$ is $\equiv \mod 2$ then $\{-36, 15\}$ is a system of representatives of $\equiv \mod 2$ and $\mathbb{Z} = \overline{-36} \lor \overline{15}$.

In later sections, this shall be very useful, so we give it the name of

Mantra 5.13. of Equivalence Relations. In the above setup, we have

$$A = \bigvee_S \overline{x}.$$  

In particular, if $|A| < \infty$, then

$$|A| = \sum_S |\overline{x}|.$$

This is only useful if we can compute the size of the equivalence class $\overline{x}$.

Exercises 5.14.

1. A partition of a set $A$ is a collection $\mathcal{C}$ of subsets of $A$ such that $A = \bigvee_{\mathcal{C}} B$. Let $R$ be an equivalence relation on $A$. Show that $\overline{A}$ partitions $A$. Conversely, let $\mathcal{C}$ partition $A$. Define a relation $\sim$ on $A$ by $a \sim b$ if $a$ and $b$ belong to the same set in $\mathcal{C}$. Show $\sim$ is an equivalence relation on $A$. So an equivalence relation and a partition of a set are essentially the same.

2. Draw lines in $\mathbb{R}^2$ perpendicular to the $X$-axis through all integer points and the analogous lines parallel to the $Y$-axis. Define a partition of the plane that arises. [Be careful about the points on the lines.]
   
   [Question. Can you do this in such a way that when viewed in $\mathbb{R}^3$ the corresponding equivalence classes looks like a torus?]

6. Modular Arithmetic

We introduce one of the most important equivalence relations in elementary arithmetic. It generalizes $\equiv \mod 2$. 
Definition 6.1. Fix an integer $m > 1$ and let $a, b \in \mathbb{Z}$. We say that $a$ is congruent to $b$ modulo $m$ and write $a \equiv b \mod m$ if $m \mid a - b$ in $\mathbb{Z}$.

The set 
$$[a]_m := \{ x \in \mathbb{Z} \mid x \equiv a \mod m \}$$

is a subset of $\mathbb{Z}$ called the residue class of $a$ modulo $m$. We shall denote this set as above or as $a + m\mathbb{Z}$. For example,
$$m\mathbb{Z} = 0 + m\mathbb{Z} = \{ km \mid k \in \mathbb{Z} \},$$

all the multiples of $m$.

Question 6.2. In the above, what would happen if we let $m = 1$ or $m = 0$? What if we let $m < 0$?

Let $a$ and $b$ be integers and $m > 1$ an integer. We have three ways of saying the same thing. They are:

$$m \mid a - b \quad \text{or} \quad a \equiv b \mod m \quad \text{or} \quad \overline{a} = \overline{b}.$$ 

One usually uses the one that is most convenient, although algebraically the last is the most interesting.

Proposition 6.3. Let $m > 1$ in $\mathbb{Z}$. Then $\equiv \mod m$ is an equivalence relation. In particular,
$$\mathbb{Z} = \overline{0} \lor \overline{1} \lor \cdots \lor \overline{m-1}.$$ 

Write $\mathbb{Z}/m\mathbb{Z}$ for $\mathbb{Z}/\equiv \mod m$. Then
$$\mathbb{Z}/m\mathbb{Z} = \{ \overline{0}, \overline{1}, \cdots, \overline{m-1} \} \text{ and } |\mathbb{Z}/m\mathbb{Z}| = m.$$ 

Let $a, b, c, d \in \mathbb{Z}$ satisfy
$$a \equiv b \mod m \text{ and } c \equiv d \mod m.$$ 

Then
$$a + c \equiv b + d \mod m \text{ and } a \cdot c \equiv b \cdot d \mod m.$$ 

Equivalently,
$$\text{if } \overline{a} = \overline{b} \text{ and } \overline{c} = \overline{d}, \text{ then } \overline{a + c} = \overline{b + d} \text{ and } \overline{a \cdot c} = \overline{b \cdot d}.$$ 

Define
$$+ : \mathbb{Z}/m\mathbb{Z} \times \mathbb{Z}/m\mathbb{Z} \rightarrow \mathbb{Z}/m\mathbb{Z} \text{ by } \overline{a} + \overline{b} := \overline{a + b}$$ 

and
$$\cdot : \mathbb{Z}/m\mathbb{Z} \times \mathbb{Z}/m\mathbb{Z} \rightarrow \mathbb{Z}/m\mathbb{Z} \text{ by } \overline{a} \cdot \overline{b} := \overline{a \cdot b}.$$ 

Then $+$ and $\cdot$ are well-defined.
[To be well-defined means that they are functions. In this case, this means that if \( a = a_1 \) and \( b = b_1 \), then \( a + b = a_1 + b_1 \) and \( a \cdot b = a_1 \cdot b_1 \).]

Moreover, \((\mathbb{Z}/m\mathbb{Z}, +, \cdot)\) satisfies the axioms of a commutative ring, i.e., for all \( a, b, c \in \mathbb{Z} \), we have

1. \( (a + b) + c = a + (b + c) \).
2. \( 0 + a = a = a + 0 \).
3. \( a + (\neg a) = 0 = (\neg a) + a \).
4. \( a \cdot b = b \cdot a \).
5. \( (a \cdot b) \cdot c = a \cdot (b \cdot c) \).
6. \( 1 \cdot a = a = a \cdot 1 \).
7. \( a \cdot (b + c) = a \cdot b + a \cdot c \).
8. \( (b + c) \cdot a = b \cdot a + c \cdot a \).

**Proof.** Exercise. \(\square\)

The element \( 0 \) is called the zero or the unity of \( \mathbb{Z}/m\mathbb{Z} \) under + and \( 1 \) is called the one or the unity of \( \mathbb{Z}/m\mathbb{Z} \) under \( \cdot \).

**Notation 6.4.** As is customary, we shall usually drop the multiplication symbol \( \cdot \) in \( a \cdot b \) when convenient.

**Remark 6.5.** Let \( \sim \) be an equivalence relation on \( A \). If \( B \) is a set, to show that \( f : A \to B \) is well-defined, i.e., \( f \) is a function, one must show that

\[
\overline{a} = \overline{a'} \implies f(\overline{a}) = f(\overline{a'}),
\]

i.e., is independent of the representative for \( \overline{a} \).

The proposition then also says the canonical surjection

\[
\pi : \mathbb{Z} \to \mathbb{Z}/m\mathbb{Z} \text{ given by } x \mapsto \overline{x} (= [x]_m)
\]

satisfies

\[
\overline{a + b} = \overline{a} + \overline{b} \text{ and } \overline{a \cdot b} = \overline{a} \cdot \overline{b}
\]

and

\[
0 \mapsto \overline{0} \text{ and } 1 \mapsto \overline{1}.
\]

**Definition 6.6.** A commutative ring is a set \( R \) together with two maps

\[
+: R \times R \to R \text{ and } \cdot : R \times R \to R,
\]

write \( +(a, b) \) as \( a+b \) and \( \cdot(a, b) \) as \( a \cdot b \), called addition and multiplication, respectively, satisfying for all \( a, b, c \in R \) the axioms (1)–(9) above hold with \( R \) replacing \( \mathbb{Z}/m\mathbb{Z} \). In (2), there exists an element 0 satisfying (2) and (3) and in (6) an element 1 satisfying (6). In particular, \( R \) is not empty. If we do not necessarily assume that \( R \) satisfies (7), we call \( R \) a ring.
Examples 6.7. 1. The sets $\mathbb{Z}$, $\mathbb{Q}$, $\mathbb{R}$, $\mathbb{C}$, $\mathbb{Z}/m\mathbb{Z}$ are commutative rings.

2. If $R$ is a ring, then so are $M_n(R)$ and $R[t]$ under the usual $+$ and $\cdot$ of matrices and polynomials, respectively.

We call a map $f : R \to S$ of rings a ring homomorphism if it preserves $+$ and $\cdot$ and takes unities to unities. For example, 

$- : \mathbb{Z} \to \mathbb{Z}/m\mathbb{Z}$ is a surjective ring homomorphism.

[A surjective ring homomorphism is also called an epimorphism.]

To investigate the interrelationship of congruences of different moduli is an appropriate study. We turn to the important and useful case of congruences of pairwise relatively prime moduli. We begin with the following lemma.

Lemma 6.8. Let $m$, $n$, and $a_i$, $1 \leq i \leq r$, be integers.

1. If $(a_i, m) = 1$ for $i = 1, \ldots, r$, then $(a_1 \cdots a_r, m) = 1$.
2. If $(a_i, a_j) = 1$ for $i \neq j$ and $a_i | n$ for $i, j = 1, \ldots, r$, then $a_1 \cdots a_r | n$.

Proof. (1): By induction, it suffices to do the case $r = 2$. (Why?) By Key Observation (4.13), we have equations 

$$x_1 a_1 + y_1 m = 1 = x_2 a_2 + y_2 m,$$

for some $x_1, x_2, y_1, y_2 \in \mathbb{Z}$, so 

$$1 = (x_1 a_1 + y_1 m)(x_2 a_2 + y_2 m) = x_1 x_2 a_1 a_2 + \text{something} \cdot m$$

and we are done.

[Note. We only used induction for the convenience of notation – do you see why?]

(2): The case $r = 1$ is immediate. By induction, we have $a_1 \cdots a_{r-1} | n$. By (1), we have $(a_1 \cdots a_{r-1}, a_r) = 1$ (setting $m = a_r$). By Key Observation (4.13), there exists an equation 

$$a_1 \cdots a_{r-1} x + a_r y = 1 \text{ for some } x, y \in \mathbb{Z}.$$ 

By Key Trick (4.14), we have 

$$a_1 \cdots a_{r-1} n x + a_r ny = n.$$ 

As $a_1 \cdots a_{r-1} a_r | a_1 \cdots a_{r-1} n$ and $a_1 \cdots a_{r-1} a_r | a_r n$, we conclude that $a_1 \cdots a_r | n$, as needed. $\square$

Theorem 6.9. (Chinese Remainder Theorem) Let $m_i$ be integers with $(m_i, m_j) = 1$ for $1 \leq i, j \leq r$ and $i \neq j$. Set $m = m_1 \cdots m_r$ and
suppose that $c_1, \ldots, c_r$ are integers. Then there exists an integer $x$ satisfying all of the following:

\[ x \equiv c_1 \mod m_1 \]
\[ \vdots \]
\[ x \equiv c_r \mod m_r. \]

Moreover, the integer $x$ is unique modulo $m$, i.e., if an integer $y$ also satisfies $y \equiv c_i \mod m_i$ for $1 \leq i \leq r$, then $x \equiv y \mod m$.

**Proof.** **Existence:** Let $n_i = \frac{m}{m_i} = m_1 \cdots \hat{m}_i \cdots m_r$, where $\hat{\cdot}$ means omit. By Lemma 6.8(1), we have $(m_i, n_i) = 1$ for $1 \leq i \leq r$, so by Key Observation (4.13), there exist equations

\[ 1 = d_im_i + e_in_i, \]

for some integers $d_i, e_i, i = 1, \ldots, r$. Set $b_i = e_in_i$ for $i = 1, \ldots, r$. Then

\[ 1 = d_im_i + b_i \text{ and } m_j | b_i \text{ if } i \neq j. \]

This means that

\[ 1 \equiv b_i \mod m_i \text{ and } 0 \equiv b_i \mod m_j \text{ if } i \neq j. \]

Consequently,

\[ x_0 := c_1b_1 + \cdots + c_rb_r \equiv c_ib_i \equiv c_i \mod m_i, \text{ with } i = 1, \ldots, r \]

and $x_0$ works.

**Uniqueness:** Suppose that $y_0$ also works. Then $x_0 \equiv y_0 \mod m_i$ for $1 \leq i \leq r$, i.e., $m_i | x_0 - y_0$ for $1 \leq i \leq r$. By Lemma 6.8 (2), we have $x_0 \equiv y_0 \mod m$. □

We want to interpret what we did above in the language of “rings” and “ring homomorphisms”. Let $m > 1$ in $\mathbb{Z}$ and $(a,m) = 1$. By Key Observation (4.13), there are integers $x$ and $y$ satisfying $1 = ax + my$. Applying the canonical surjection $\pi : \mathbb{Z} \to \mathbb{Z}/m\mathbb{Z}$ to this yields

\[ \overline{1} = ax + my = \overline{a} \overline{x} + \overline{m} \overline{y} = \overline{a} \overline{x} + \overline{0} \overline{y} = \overline{a} \overline{x} (= \overline{x \overline{a}}), \]

i.e., $\overline{a}$ has a multiplicative inverse in $\mathbb{Z}/m\mathbb{Z}$. An element $\overline{a}$ that has a multiplicative inverse is called a **unit**. Let

\[ (\mathbb{Z}/m\mathbb{Z})^\times := \{ \overline{a} \in \mathbb{Z}/m\mathbb{Z} \mid \overline{a} \text{ is a unit} \}. \]

So if $a \in \mathbb{Z}$ and $(a,m) = 1$, then $\overline{a} \in (\mathbb{Z}/m\mathbb{Z})^\times$.

Conversely, suppose $a \in \mathbb{Z}$ satisfies $\overline{a} \in (\mathbb{Z}/m\mathbb{Z})^\times$. Then there exists $b \in \mathbb{Z}$ satisfying $\overline{a} \overline{b} = \overline{1}$. In particular, we have $m | ab - 1$, so $ab - 1 = mk$ for some integer $k$, or equivalently, $ab - mk = 1$. We therefore have
Conclusion 6.10. Let $a$ and $m > 1$ be integers. Then
\[ \overline{a} \in (\mathbb{Z}/m\mathbb{Z})^\times \text{ if and only if } (a, m) = 1. \]

We can now interpret Lemma 6.8(1) and Properties 4.9(1) as follows: If $x$ and $y$ are integers then
\[ x, y \in (\mathbb{Z}/m\mathbb{Z})^\times \text{ if and only if } xy \in (\mathbb{Z}/m\mathbb{Z})^\times. \]

In particular, the set of units $(\mathbb{Z}/m\mathbb{Z})^\times$ of $\mathbb{Z}/m\mathbb{Z}$ is closed under multiplication (but not addition), i.e., we can write
\[ \cdot : (\mathbb{Z}/m\mathbb{Z})^\times \times (\mathbb{Z}/m\mathbb{Z})^\times \rightarrow (\mathbb{Z}/m\mathbb{Z})^\times \text{ taking } (\overline{a}, \overline{b}) \rightarrow \overline{a} \cdot \overline{b}. \]

We wish to generalize this.

Definition 6.11. Let $R$ be a ring with $1 \neq 0$ and set
\[ R^\times := \{ r \in R \mid \text{there exists an } a \in R \text{ such that } ar = 1 = ra \}, \]
the units of the ring $R$. Note that this set is closed under multiplication.

Question 6.12. What are $\mathbb{Z}^\times$, $\mathbb{Q}^\times$, $\mathbb{R}^\times$, $\mathbb{C}^\times$, $\mathbb{M}_n(\mathbb{R})^\times$, $\mathbb{M}_n(\mathbb{Z})^\times$?

We now interpret the Chinese Remainder Theorem in this new language. As before, let $m_i$, $1 \leq i \leq r$, be integers with $(m_i, m_j) = 1$ if $i \neq j$, and $m = m_1 \cdots m_r$. Then $m_j \mid m$ for all $j$. In particular, if $m \mid a - b$, then $m_j \mid a - b$ for all $j$, i.e., if $a \equiv b \mod m$, then $a \equiv b \mod m_j$ for all $j$. This means that the map
\[ Z/m\mathbb{Z} \rightarrow Z/m_j\mathbb{Z} \text{ given by } [a]_m \mapsto [a]_{m_j} \]
is well-defined (i.e., a function). For example, we have $7 \equiv 1 \mod 6$, so $7 \equiv 1 \mod 3$ and $7 \equiv 1 \mod 2$. This leads to a map
\[ Z/m\mathbb{Z} \rightarrow Z/m_1\mathbb{Z} \times \cdots \times Z/m_r\mathbb{Z} \]
(6.13)
given by $[a]_m \mapsto ([a]_{m_1}, \ldots, [a]_{m_r})$.

Giving the right hand set component-wise operations turns it into a commutative ring and we see immediately that this map is a ring homomorphism. (Note that the zero of the right hand side is $([0]_{m_1}, \ldots, [0]_{m_r})$ and the one is $([1]_{m_1}, \ldots, [1]_{m_r})$.) The Chinese Remainder Theorem says that this map is bijective. The inverse of this map is also checked to be a ring homomorphism. So the Chinese Remainder Theorem says the map above is a ring isomorphism, i.e., a bijective ring homomorphism whose inverse is also a ring homomorphism. (We shall see that this last condition is unnecessary.) This means that the rings
\[ Z/m\mathbb{Z} \text{ and } Z/m_1\mathbb{Z} \times \cdots \times Z/m_r\mathbb{Z} \]
are isomorphic, i.e., they look the same algebraically. For example, if $m$ and $n$ are relatively prime integers greater than 1, then $Z/mn\mathbb{Z}$ and $Z/m\mathbb{Z} \times$
Z/nZ are isomorphic. In particular, if \( n = p_1^{e_1} \cdots p_s^{e_s} \) is a standard factorization, then \( \mathbb{Z}/n\mathbb{Z} \) and \( \mathbb{Z}/p_1^{e_1}\mathbb{Z} \times \cdots \times \mathbb{Z}/p_s^{e_s}\mathbb{Z} \) are isomorphic. This last fact is useful as it reduces solving many equations modulo \( n \) to congruences modulo the various prime power constituents in its standard factorization.

It is left as an exercise to show the isomorphism in (6.13) above induces a bijection between

\[
(\mathbb{Z}/m\mathbb{Z})^\times \text{ and } (\mathbb{Z}/m_1\mathbb{Z})^\times \times \cdots \times (\mathbb{Z}/m_r\mathbb{Z})^\times.
\]

This means if we define the Euler phi-function \( \varphi : \mathbb{Z}^+ \to \mathbb{Z}^+ \) by \( \varphi(1) = 1 \) and \( \varphi(m) = |(\mathbb{Z}/m\mathbb{Z})^\times| \) for \( m > 1 \), then \( \varphi(mn) = \varphi(m)\varphi(n) \) whenever \( m \) and \( n \) are relatively prime positive integers.

**Exercises 6.14.**

1. (Fermat's Little Theorem) Let \( p \) be a prime. Then for all integers \( a \), we have

\[
a^p \equiv a \mod p.
\]

[You can use the Binomial Theorem if you state it carefully and then prove the *Children's Binomial Theorem* that says for all integers \( a \) and \( b \), we have \( (a + b)^p \equiv a^p + b^p \mod p \). We will have another proof of this later.]

2. Prove that there exist infinitely many primes \( p \) satisfying \( p \equiv 3 \mod 4 \). [It is true that there exist infinitely many primes \( p \) satisfying \( p \equiv 1 \mod 4 \) but this is harder.]

[Hint: Look at the proof for the infinitude of primes.]

3. Prove Proposition 6.3.

4. Find the smallest positive integer \( x \) satisfying the congruences

\[
x \equiv 3 \mod 11 \quad x \equiv 2 \mod 12 \quad \text{and} \quad x \equiv 3 \mod 13.
\]

5. What are the units of the rings \( \mathbb{Q}, \mathbb{Z}, \mathbb{R}[t], \mathbb{Z}[t], (\mathbb{Z}/4\mathbb{Z})[t], \) and \( \mathbb{M}_n(\mathbb{C}) \)?

6. Show if \( m \) and \( n \) are relatively prime positive integers, then the map

\[
(\mathbb{Z}/mn\mathbb{Z})^\times \to (\mathbb{Z}/m\mathbb{Z})^\times \times (\mathbb{Z}/n\mathbb{Z})^\times
\]

given by \([a]_{mn} \mapsto ([a]_m, [a]_n)\) is a bijection. In particular, \( \varphi(mn) = \varphi(m)\varphi(n) \) whenever \( m \) and \( n \) are relatively prime positive integers.

7. Let \( n \) and \( r \) be positive integers and \( p \) a positive prime. Show that

\[
\varphi(p^r) = p^{r-1}(p-1) \quad \text{and} \quad \varphi(n) = n \prod_{\substack{p|n \\text{ prime}}}(1 - \frac{1}{p}).
\]
7. Surjective maps

In the previous sections, we introduced equivalence relations, using them to construct the integers mod \( m \) and the corresponding surjection from the integers to the integers mod \( n \). We want to expand the usefulness of equivalence relations. In mathematics, one tries to reduce the study of functions that preserve a given structure into injective and surjective maps that preserve that structure. The reason for studying injective maps is that if one can embed, i.e., find an injection of an object into one with more structure, one can often pull back information. The reason to study surjective maps is that usually the target is simpler, e.g., \( \mathbb{Z}/m\mathbb{Z} \) is simpler than \( \mathbb{Z} \). If we are just studying sets, i.e., without any additional structure, this reduction can always be accomplished. Equivalence relations give the surjective maps.

Exercise 5.14(1) showed that equivalence relations are essentially the same as partitioning. Let \( A \) be a set with an equivalence relation \( \sim \) on it. Then the set of equivalence classes, \( \overline{A} := \{ \overline{a} \mid a \in A \} \), partitions \( A \), and the exercise says that partitions of \( A \) induce equivalence relations on \( A \). Let \( S \) be a system of representatives for \( A \) relative to \( \sim \), so \( A = \bigcup_S \overline{a} \). We then have the canonical surjection

\[ \pi : A \to \overline{A} \text{ given by } a \mapsto \overline{a}, \]

i.e., every equivalence relation leads to a surjective map. We want the converse. Let \( f : A \to B \) be a surjective map.

We define an equivalence relation \( \sim \) on \( A \) by

\[ a \sim a' \text{ if } f(a) = f(a'). \]

This is clearly an equivalence relation on \( A \), as \( = \) is an equivalence relation on \( B \). We compute the equivalence class of \( a \) to be

\[ \overline{a} = \{ x \in A \mid x \sim a \} = \{ x \in A \mid f(x) = f(a) \}, \]

called the fiber of \( f \) at \( f(a) \), and let

\[ f^{-1}(f(a)) = \{ x \in A \mid f(x) = f(a) \}. \]

[In general, \( f^{-1} \) is not a function, but we use the general notation if \( f : A \to B \) is any map and \( D \subset B \), then \( f^{-1}(D) := \{ x \in A \mid f(x) \in D \} \). If \( D = \{ b \} \) (the preimage of \( D \)), then the fiber \( f^{-1}(b) \) at \( b \) is \( f^{-1}(\{ b \}) \).]

We now know, if \( \sim \) is the equivalence relation defined by the surjective map \( f \), then we have

\[ (*) \quad \overline{a} = \overline{a'} \text{ if and only if } a \sim a' \text{ if and only if } f(a) = f(a'). \]
But this means if we define
\[ \overline{f} : \overline{A} \to B \text{ by } \overline{a} \mapsto f(a), \]
then
\[ \overline{f}(\overline{a}) = f(a), \]
so \( \overline{f} \) is a well-defined injective map by (*). We say that \( f \) induces \( \overline{f} \). This induced map is also surjective since \( \overline{f}(\overline{A}) = f(A) \), hence it is a bijection. So set theoretically we cannot tell the sets \( \overline{f}(\overline{A}) \) and \( B \) apart. Moreover, we have a commutative diagram
\[
\begin{array}{ccc}
A & \xrightarrow{f} & B \\
\downarrow & & \uparrow \\
\overline{A} & \xrightarrow{\overline{f}} & \overline{B}
\end{array}
\]
i.e., \( f = \overline{f} \circ \overline{\cdot} \). We call this the First Isomorphism Theorem of Sets. [In general, we say a diagram commutes if following any composition of any maps (arrows) from the same place to the same target gives equal maps.]

This is so useful that we summarize the above.

**Summary 7.1.** Let \( f : A \to B \) be a surjective map. Define an equivalence relation \( \sim \) on \( A \) by \( a \sim a' \) if \( f(a) = f(a') \). Then

1. \( \sim \) is an equivalence relation.
2. \( \overline{A} = \{ f^{-1}(f(a)) \mid a \in A \} \).
3. \( \overline{\cdot} : A \to \overline{A} \) given by \( a \) is mapped to \( \overline{a} \) is a surjective map.
4. \( \overline{f} : \overline{A} \to B \) given by \( \overline{a} \) maps to \( f(a) \) is a well-defined bijection.
5. The diagram

\[
\begin{array}{ccc}
A & \xrightarrow{f} & B \\
\downarrow & & \uparrow \\
\overline{A} & \xrightarrow{\overline{f}} & \overline{B}
\end{array}
\]
commutes.

We can push this a bit further. Let \( g : A \to C \) be a map and set \( B = g(A) \). Let \( f : A \to B \) be given by \( f(a) = g(a) \) for all \( a \in A \), i.e., we change the target. Then \( g \) is the composition

\[
\begin{array}{ccc}
A & \xrightarrow{f} & B & \xrightarrow{\text{inc}} & C
\end{array}
\]
with inc the inclusion map.

**Notation 7.2.** In the future, we shall usually abuse notation and write the same letter for a function and the function arising by changing its target (if it makes sense) when no confusion can arise.
For clarity, we still write \( g \) for this \( f \) here. Let \( \sim \) be the equivalence relation given by \( a \sim a' \) if \( f(a) = f(a') \). Then we have a commutative diagram

\[
\begin{array}{ccc}
A & \xrightarrow{g} & C \\
\downarrow & \searrow & \nearrow \text{inc} \\
\bar{A} & \xrightarrow{\bar{f}} & B.
\end{array}
\]

If we set \( \bar{g} = \text{inc} \circ \bar{f} \), we get a commutative diagram

\[
\begin{array}{ccc}
A & \xrightarrow{g} & C \\
\downarrow & \searrow & \nearrow \bar{g} \\
\bar{A} & \xrightarrow{\bar{f}} & B.
\end{array}
\]

with \( \sim \) a surjection and \( \bar{g} \) an injection. This is the First Isomorphism Theorem (alternate version), showing that every map factors as a composition injection \( \circ \) surjection. Note that the equivalence classes of \( \sim \) are precisely the non-empty fibers of the map \( g \).

If we give our sets additional structure, we would be interested in maps that preserve this additional structure. In general, the map \( \bar{g} \) may not preserve the additional structure, and the set \( \bar{A} \) may not have a compatible structure. We shall see in the sequel that in certain cases we can achieve this.

The map \( g : A \to C \) is also a composition of surjection \( \circ \) injection. Indeed let \( \Gamma_g : A \to A \times C \) be given by \( a \mapsto (a, g(a)) \) (called the graph of \( g \)) and \( \pi_C : A \times C \to C \) by \( (a, c) \mapsto c \), the projection of \( A \times C \) onto \( C \). Then \( g = \pi \circ \Gamma_g \). This shows that to study (set) maps, it suffices to study injections and surjections.
Part 2

Group Theory
CHAPTER III

Groups

In this chapter, we begin our study of abstract algebra. The basic object in this study is a group, a set \( G \) with one (binary) operation \( \circ \) subject to three axioms: associativity, existence of a unity, and existence of inverses. An example that you know from linear algebra is the set \( \text{GL}_n(\mathbb{R}) \) of all \( n \times n \) matrices over the reals with nonzero determinant which is called the \( n \times n \) general linear group over \( \mathbb{R} \). Although groups are easy to define, their structure is far from simple. A primary reason for this is that, in general, we do not assume that \( x \circ y = y \circ x \) in \( G \).

In algebra, as in other fields of mathematics, one studies sets with additional structure by investigating maps between those sets that preserve the additional structure. In particular, one would like to know when two objects with additional structure are essentially the same. Here the prototypes that you have encountered in linear algebra are linear transformations between vector spaces and isomorphisms of vector spaces. One idea to keep in mind is that in modern mathematics maps between objects are more important than the objects themselves. In this chapter, we give many examples of groups and maps between them and establish the basic theorems needed to investigate groups.

It is important that you learn the many examples given below (as well as others), as you can only theorize what is true based on examples that you know – of course, your guess may very well not be true. In this course, we shall mostly apply our theorems to groups having finitely many elements because this allows us to use the power of equivalence relations in an effective way – we can count. The general theory of arbitrary groups is very difficult, but it is extremely important not only in algebra but in other fields of study because of the general linear group (and its subgroups).

8. Definitions and Examples

A map \( \cdot : G \times G \to G \) is called a binary operation. We shall always write \( a \cdot b \) for \( \cdot(a, b) \).
Definition 8.1. Let $G$ be a set with a binary operation $\cdot : G \times G \to G$. We call $(G, \cdot)$ a group if it satisfies the following axioms:

**Associativity.** $(a \cdot b) \cdot c = a \cdot (b \cdot c)$ for all $a, b, c \in G$.

**Unity.** There exists an element $e \in G$ such that for all $a \in G$, we have $a \cdot e = a = e \cdot a$. The element $e$ is called a unity or an identity.

**Inverses.** There exists a unity $e$ in $G$; and, for any $x \in G$, there exists a $y \in G$ satisfying $xy = e = yx$.

If, in addition, $G$ satisfies

**Commutativity.** $a \cdot b = b \cdot a$ for all $a, b \in G$,

we call $G$ an abelian group.

We usually write $G$ for $(G, \cdot)$ and $ab$ for $a \cdot b$ if $\cdot$ is clear.

The most common algebraic property on a set with a binary operation is associativity, although there are interesting algebraic objects that do not satisfy this property. However, all the algebraic objects that we shall study will satisfy an analogue of associativity. If our set $G$ under its binary operation only satisfies Associativity and Unity, then it is called a monoid.

Remarks 8.2. Let $\cdot : G \times G \to G$ be a binary operation.

1. If $G$ satisfies associativity and $a_1, \ldots, a_n \in G$, then $a_1 \cdots a_n$ makes sense, i.e., is independent of parentheses. In particular, if $n \in \mathbb{Z}^+$ then $a^n$ makes sense with $a^1 = a$ and $a^n = a(a^{n-1})$ for $n > 1$. If $G$ is a monoid, we let $a^0 = e$. The nasty uninteresting induction on the complexity of parentheses is left to the diligent reader.

2. If $G$ satisfies Unity, then the unity $e$ is unique. Indeed if $e'$ is another unity then $e = ee' = e'$.

[The unity of an abstract group will usually be written $e_G$ or $e$ but for many specific types of groups the unity will be written as 0 or 1.]

3. If $G$ is a monoid, then $a \in G$ has at most one inverse. If it has an inverse, then it is denoted by $a^{-1}$, in which case $a$ is the inverse of $a^{-1}$. Indeed if $b$ and $c$ are inverses of $a$ then

$$b = b \cdot e = b \cdot (a \cdot c) = (b \cdot a) \cdot c = e \cdot c = c.$$  

[In fact, the proof shows if there exist $b, c \in G$ satisfying $ba = e = ac$ then $b = c$ is the inverse of $a$.]

4. If $G$ is a monoid and $a, b \in G$ have inverses then so does $ab$ and $(ab)^{-1} = b^{-1}a^{-1}$. In addition, if $n \in \mathbb{Z}^+$ then $a^{-n} = (a^{-1})^n$.

**Question.** Let $a, b$ be elements in a monoid such that $ab$ has an inverse. Is it true that $a$ and $b$ have inverses?
5. If \( G \) is a group, then the **cancellation laws** hold, i.e., for all \( a, b, c \in G \), we have

\[
ab = ac \implies b = c \quad \text{and} \quad ba = ca \implies b = c.
\]

In particular, if \( a = xy \) in \( G \), then \( x = ay^{-1} \) and \( y = x^{-1}a \).

6. If \( G \) is a group under a binary operation notated by \( + : G \times G \to G \), then \( G \) will always be an abelian group. We then call \( G \) an **additive group** and write \( 0 \) (or \( 0_G \)) for \( e_G \) and \( -a \) for \( a^{-1} \).

**Definition 8.3.** Using this new language, it is easy to define a **ring** more carefully. Let \( R \) be a set with two binary operations \( \cdot : R \times R \to R \) and \( + : R \times R \to R \). Then \( R \) is a ring under **addition** \( + \) and **multiplication** \( \cdot \) if \( (R, +) \) is an additive group, \( (R, \cdot) \) is a monoid, and \( R \) satisfies the **Distributive Laws**

\[
a \cdot (b + c) = a \cdot b + a \cdot c \quad \text{and} \quad (b + c) \cdot a = b \cdot a + c \cdot a,
\]

for all \( a, b, c \in R \). The distributive laws interrelate the two binary operations. The multiplicative unity is written \( 1 \) (or \( 1_R \) if there can be confusion). The ring is called a **commutative ring** if \( (R, \cdot) \) is a commutative monoid. One checks that \( 1 = 0 \) in \( R \) if and only if \( R = \{0\} \), the zero or **trivial ring**. If \( R \) is not the zero ring it is called a **division ring** if \( (R \setminus \{0\}, \cdot) \) is a group. A commutative division ring is called a **field**.

It is worth mentioning again that when one has defined a new concept, one should always try to find many examples. Theorems usually arise from knowledge of explicit examples. We present many examples of groups.

**Examples 8.4.** 1. A **trivial** group is a group consisting of a single element, necessarily the unity.

2. \( \mathbb{Z}, \mathbb{Q}, \mathbb{R}, \mathbb{C}, \mathbb{Z}/m\mathbb{Z} \) \((m > 1)\), or any ring is an additive group under \(+\).

3. \( \mathbb{R}^+ \), the set of positive real numbers, is an abelian group but \( \mathbb{Z}^+ \) is only an abelian (i.e., commutative) monoid under multiplication with unity \( e = 1 \), and not a group. Neither are monoids under addition, although they still satisfy associativity.

4. If \( F \) is a field, e.g., if \( F = \mathbb{Q}, \mathbb{R}, \) or \( \mathbb{C} \), then by definition \( F^\times = F \setminus \{0\} \) is an abelian group under multiplication, e.g., if \( F = \mathbb{Q}, \mathbb{R}, \) or \( \mathbb{C} \). More generally, if \( R \) is any ring then its set of units \( R^\times \), i.e., the set of elements having a multiplicative inverse, is a group under \( \cdot \), abelian if \( R \) is commutative, called the **group of units of** \( R \). This leads to many examples of groups.

5. Let \( V \) be a vector space over a field \( F \). Then \( (V, +) \) is an additive group.
6. Let \( S \) be a non-empty set. Then

\[
\Sigma(S) := \{ f : S \to S \mid f \text{ is a bijection}\}
\]

is a group under the composition of functions. The unity is the identity map on \( S \) which we shall write as \( 1_S \). A bijection \( f : S \to S \) is called a permutation. [If \( f \in \Sigma(S) \), what is \( f^{-1} \)?] We call \( \Sigma(S) \) the group of all permutations of \( S \). It is also a transitive group on \( S \). That is, for all \( x, y \in S \), there exists a permutation \( f \in \Sigma(S) \) satisfying \( f(x) = y \). If \( S = \{1, \ldots, n\} \) then \( \Sigma(S) \) is denoted \( S_n \) and called the symmetric group on \( n \) letters. Note that \( |S_n| = n! \).

**Question.** Let \( S = \{a, b, \ldots, z\} \). Can you tell the difference between \( \Sigma(S) \) and \( S_{26} \) algebraically?

For the next examples, we need the definition of a subgroup of a group. A subset \( H \) of a group \( G \) is called a subgroup of \( G \) if it becomes a group under the restriction of the binary operation on \( G \) to \( H \), i.e., \( \cdot|_{H \times H} : H \times H \to H \) makes sense (meaning that the image of \( \cdot|_{H \times H} \) lies in \( H \)), and has the same unity as \( G \) (although this last condition will be seen to be unnecessary).

7. Let \( S \) be a non-empty set and \( x_0 \in S \). The set

\[
\Sigma(S)_{x_0} := \{ f \in \Sigma(S) \mid f(x_0) = x_0\} \subset \Sigma(S)
\]

is a group called the stabilizer of \( x_0 \) in \( \Sigma(S) \). We say elements of \( \Sigma(S)_{x_0} \) fix \( x_0 \). We shall see that groups often arise as functions acting on sets and a major problem is to find fixed points of these actions, e.g., in the above \( x_0 \) is a fixed point of the action of \( \Sigma(S)_{x_0} \) on \( S \). Note that \( (S_n)_n \) looks algebraically like \( S_{n-1} \).

\( \Sigma(S)_{x_0} \) is also a subgroup of \( \Sigma(S) \).

We can generalize the stabilizer of \( x_0 \) as follows: Let \( x_0, \ldots, x_n \) be elements of \( S \). Then

\[
\Sigma(S)_{x_0} \cap \cdots \cap \Sigma(S)_{x_n} = \{ f \in \Sigma(S) \mid f(x_i) = x_i, \text{ for } i = 1, \ldots, n\}
\]

is a subgroup of \( \Sigma(S) \) (as is \( \Sigma(S)_{x_i} \) for all \( i \)) stabilizing (fixing) \( x_0, \ldots, x_n \).

8. Let \( G \) be a group, \( H_i \) with \( i \in I \) be subgroups of \( G \). Then \( \bigcap_I H_i \) is a subgroup of \( G \). In general, \( \bigcup_I H_i \) is not a subgroup of \( G \). [Can you find a condition when it is?]

9. Let \( G \) be a group and \( W \subset G \) a subset. Set

\[
\mathcal{W} = \{ H \subset G \mid H \text{ is a subgroup of } G \text{ with } W \subset H\}.
\]
As $W \subset G$, we have $W \neq \emptyset$. Let 

$$\langle W \rangle := \bigcap_{H \text{ subgroup of } G} H$$

This is the unique smallest subgroup of $G$ containing $W$. [Can you show this?] We say that $W$ generates $\langle W \rangle$ and that $W$ is a set of generators for $\langle W \rangle$. A set of generators is not unique. If $W = \{a_1, \ldots, a_n\}$ is finite, we also write $\langle a_1, \ldots, a_n \rangle$ for $\langle \{a_1, \ldots, a_n\} \rangle$.

We say $G$ is finitely generated if there exists a finite set $W$ such that $G = \langle W \rangle$ and cyclic if there exists an element $a \in G$ such that $G = \langle a \rangle$. If this is the case then $G = \{a^n \mid n \in \mathbb{Z}\}$ and is abelian. For example, $(\mathbb{Z}, +)$ and $(\mathbb{Z}/m\mathbb{Z}, +), m > 1,$ are cyclic with generators 1 and $\mathbb{1}$, respectively.

**Question.** What are the analogues of the above for vector spaces? The analogy will not go too far, as the concept of linear independence almost always fails.

**Warning.** If $G$ is a group that can be generated by two elements, $G$ need not be abelian. As an example consider a figure eight in the plane and take the group generated by the two counterclockwise rotations of $360^\circ$, each one rotating around its respective circle starting from the intersection point. [Can you see why this group is not abelian?]

10. Let 

$$T := \{z \in \mathbb{C} \mid |z| = 1\},$$

where $|z| = \sqrt{z\overline{z}}$ with $\overline{z}$ the complex conjugate of $z$. This is an abelian group under multiplication, called the circle group. It is a subgroup of $\mathbb{C}^\times$. If $n \in \mathbb{Z}^+$ then 

$$\mu_n := \{z \in T \mid z^n = 1\} = \langle e^{2\pi\sqrt{-1}/n} \rangle$$

is a cyclic subgroup of $T$ called the group of $n$th roots of unity. Another subgroup of $T$ is 

$$\{z \in T \mid z \in \mu_n \text{ for some } n \in \mathbb{Z}^+\} = \bigcup_{n \in \mathbb{Z}^+} \mu_n.$$ 

Note that a subgroup of an abelian group is abelian.

11. The symmetries of a geometric object often form a group. We look at some special cases. Consider an equilateral triangle in the plane with a side on the $X$-axis labeling the ordered vertices $A, B, C$, with line segment $AB$ the base. Let $r$ be a counterclockwise rotation of $120^\circ$. The triangle looks the same but the vertices are now ordered
III. GROUPS

C, A, B, with base CA. Composing rotations shows $r^2 = r \circ r$ takes the triangle to $B, C, A$, and $r^3$ back to $A, B, C$, i.e., we have the relation $r^3$ is the identity. So $\langle r \rangle$ is a cyclic group of three elements. [Should we be able to compare this group to $\mu_3$?] Now, viewing the plane in $\mathbb{R}^3$, let $f$ denote the flip along the line of the apex of the triangle perpendicular to the base. It takes the ordered vertices $A, B, C$ to $B, A, C$. Under composition, we have the relation $f^2$ is the identity which we write as 1 and $\langle f \rangle$ is a cyclic group of two elements. [Should we be able to compare this group to $\mu_2$ and to $(\mathbb{Z}/2\mathbb{Z}, +)$ and to $(\mathbb{Z}^x, \cdot)$?] Composing $f$ and $r$ leads to the relations $f^{-1} \circ r \circ f = r^2 = r^{-1}$ and a non-abelian group with six elements, viz.,

$$\{1, r, r^2, f, fr, rf\}.$$ 

[Show this.] We say that $r$ and $f$ generate this group subject to the relations $r^3 = 1$, $f^2 = 1$, and $f^{-1}rf = r^{-1}$. Note that it follows that $rf = fr^2$ and $fr = r^2f$. It is customary to write this as

$$\langle r, f \mid r^3 = 1, f^2 = 1, f^{-1}rf = r^{-1}\rangle,$$

i.e., in the form

$$\langle \text{generators} \mid \text{relations} \rangle.$$

This group is called the symmetries of an equilateral triangle or the dihedral group of order six and denoted by $D_3$. [Compare $D_3$ with $S_3$.] More generally, consider a regular $n$-gon, $n \geq 3$, with $r$ a counterclockwise rotation of $2\pi/n$ radians and $f$ a flip along the perpendicular at the bisection point of the base. Then under composition, we get a non-abelian group with $2n$ elements. It is a group that is defined by two generators $r$ and $f$ satisfying the three relations

$$r^n = 1 \quad f^2 = 1 \quad f^{-1} \circ r \circ f = r^{-1}.$$ 

It follows that $f^{-1} \circ r \circ f = r^{n-1}$. It is called the symmetries of the regular $n$-gon or the dihedral group of order $2n$ and denoted by $D_n$. As above, this is usually written

$$D_n := \langle r, f \mid r^n = 1, f^2 = 1, f^{-1}rf = r^{-1}\rangle.$$

Note if $n > 3$ then $D_n$ and $S_n$ have a different number of elements.

12. Let $Q = \{1, -1, i, -i, j, -j, k, -k\}$. This becomes a non-abelian group called the quaternion group if we invoke the relations $(-1)^2 = 1$, $k = ij = -ji$ and $i^2 = j^2 = -1$.

13. Let $F = \mathbb{Q}, \mathbb{R}, \mathbb{C}$, or any field. Then

$$\text{GL}_n(F) := \{A \in \mathbb{M}_n(F) \mid \det A \neq 0\}$$
is a group under matrix multiplication called the *general linear group of degree* \( n \). It and its subgroups are probably the most important groups in mathematics. If \( n = 1 \) then \( \text{GL}_1(F) = F^\times \) but if \( n > 1 \), this group is not abelian. [Proof?]

More generally, let \( R \) be any ring. Then the set of units \((\mathbb{M}_n(R))^\times\) is a group under multiplication of matrices. If \( R \) is commutative, it turns out that the determinant of a matrix still makes sense and \( A \in (\mathbb{M}_n(R))^\times \) if and only if \( \det(A) \in R^\times \), so we let

\[
\text{GL}_n(R) := (\mathbb{M}_n(R))^\times
\]

also called the *general linear group of degree* \( n \). If \( A \in \mathbb{M}_n(R) \), we denote its \( ij \)th entry by \( A_{ij} \). Then as usual, its *transpose* is \( A^t \), where \((A^t)_{ij} = A_{ji}\) and if \( R = \mathbb{C} \), its *adjoint* or complex conjugate transpose is \( A^* \) where \((A^*)_{ij} = \overline{A_{ji}}\). If \( A \) is a diagonal matrix, we also write \( A = \text{diag}(a_1, \ldots, a_n) \), e.g., the identity is the matrix \( I = I_n = \text{diag}(1, \ldots, 1) \).

We define some of the interesting subgroups of the general linear group over a commutative ring \( R \) (the same definitions hold for over an arbitrary ring except when the determinant is involved).

- \( \text{SL}_n(R) = \{ A \in \text{GL}_n(R) \mid \det A = 1 \} \) *special linear group*
- \( \text{O}_n(R) = \{ A \in \text{GL}_n(R) \mid A^t A = I \} \) *orthogonal group*
- \( \text{SO}_n(R) = \text{SL}_n(R) \cap \text{O}_n(R) \) *special orthogonal group*
- \( \text{D}_n(R) = \{ A \in \text{GL}_n(R) \mid A \text{ diagonal} \} \) *diagonal group*
- \( \text{T}_n(R) = \{ A \in \text{GL}_n(R) \mid A_{ij} = 0 \text{ if } i > j \} \) *upper triangular group*
- \( \text{ST}_n(R) = \{ A \in \text{T}_n(R) \mid A_{ii} = 1 \text{ all } i \} \) *strictly upper triangular group*

Let \( J = \begin{pmatrix} 0 & I_n \\ -I_n & 0 \end{pmatrix} \) in \( \text{GL}_{2n}(R) \). Then

\[ \text{Sp}_{2n}(F) = \{ A \in \text{GL}_n(R) \mid A^t J A = J \} \] *symplectic group*

Let \( I_{3,1} = \text{diag}(1,1,1,-1) \) in \( \text{GL}_4(\mathbb{R}) \). Then

\[ \text{O}_{3,1}(\mathbb{R}) = \{ A \in \text{GL}_4(\mathbb{R}) \mid A^t I_{3,1} A = I_{3,1} \} \] *Lorenz Group*

\[ \text{U}_n(\mathbb{C}) = \{ A \in \text{GL}_n(\mathbb{C}) \mid A^* A = I \} \] *unitary group*

\[ \text{SU}_n(\mathbb{C}) = \text{SL}_n(\mathbb{C}) \cap \text{U}_n(\mathbb{C}) \] *special unitary group.*

14. Let \( V \) be a vector space over a field \( F \). Then

\[ \text{Aut}_F(V) := \{ T : V \to V \mid T \text{ a linear (i.e., vector space) isomorphism} \} \]
is a group under composition called the automorphism group of $V$. An isomorphism of a vector space to itself is called an automorphism. This generalizes greatly, and its generalization probably is the most important type of group, e.g., cf. $\text{Aut}_F(V)$ and $\text{GL}_n(F)$ if $V$ is an $n$-dimensional vector space.

15. Let $G_i, i \in I$, be groups and $G = \times_i G_i$ the cartesian product of the sets $G_i, i \in I$. We write elements in $G$ by $(g_i)_I$. [Technically, $\times_i G_i := \{ f : I \rightarrow \bigcup_i G_i \mid f(i) \in G_i \text{ for all } i \in I \}$.]

Then $G$ is a group under component-wise operation, e.g., $e_G = (e_{G_i})_I$. This group is called the external direct product of the $G_i$’s. If all the $G_i$ are abelian, so is $\times_i G_i$.

For example, writing $\mathbb{Z}/2\mathbb{Z} = \{0, 1\}$, we have an abelian group

$$V := \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2\mathbb{Z} = \{(0, 0), (0, 1), (1, 0), (1, 1)\}.$$ 

It satisfies:

- $(a, b) + (a, b) = (0, 0)$ if $a, b \in \mathbb{Z}/2\mathbb{Z}$
- $(1, 0) + (0, 1) = (1, 1)$
- $(1, 1) + (0, 1) = (1, 0)$
- $(1, 1) + (1, 0) = (0, 1)$

etc.

This group is called the Klein Four (Vier) Group. Note that this is an abelian group that is not cyclic. Also it is worth noting that $V$ is a two-dimensional vector space over the field $\mathbb{Z}/2\mathbb{Z}$.

16. Let $m > 0$. Then we know that $(\mathbb{Z}/m\mathbb{Z}, +)$ is cyclic and $((\mathbb{Z}/m\mathbb{Z})^\times, \cdot)$ is abelian but the latter may not be cyclic. What is $(\mathbb{Z}/8\mathbb{Z})^\times$ or more generally $(\mathbb{Z}/2^n\mathbb{Z})^\times$ if $n > 2$? The answer is interesting, as it caused a well-known theorem in number theory called Grünwald’s Theorem to be in fact false. The correction by Wang had to take the answer to this question into account.

17. If $a, b \in \mathbb{Z}^+$ have $d$ as their gcd, then the additive group $\langle a, b \rangle = \langle d \rangle$.

Exercises 8.5.

1. Let $a, b$ be elements in a monoid such that $ab$ has an inverse. Is it true that $a$ and $b$ have inverses? Prove this if true or give a counterexample if false.
2. If \( H_i, i \in I, \) are subgroups of a group \( G, \) show that \( \bigcap_I H_i \) is a subgroup of \( G, \) but \( \bigcup_I H_i \) is not a subgroup in general (i.e., give a counterexample). Find a nontrivial condition on the \( H_i \) such that \( \bigcup_I H_i \) is a subgroup.

3. Show if \( G \) is a group in which \((ab)^2 = a^2b^2\) for all \( a, b \in G, \) then \( G \) is abelian.

4. Determine all groups having at most six elements.

5. Show that the quaternion group is a group of eight elements.

6. Show that the dihedral group \( D_n \) has \( 2n \) elements.

7. Let \( G \) be a group and \( H \) a subgroup of \( G. \) Define the core of \( H \) in \( G \) by \( \text{Core}_G(H) := \bigcap_G xHx^{-1}. \) Show \( \text{Core}_G(H) \) is the unique largest subgroup \( N \) of \( H \) satisfying \( xNx^{-1} = N \) for all \( x \) in \( G. \)

8. Let \( G \) be a group and \( H \) a subgroup of \( G. \) Define the normalizer of \( H \) in \( G \) by \( \text{N}_G(H) := \{ x \in G \mid xHx^{-1} = H \}. \) Show that \( \text{N}_G(H) \) is a subgroup of \( G \) containing \( H \) and is the unique largest subgroup \( K \) of \( G \) containing \( H \) satisfying \( H = xHx^{-1} \) for all \( x \) in \( K. \)

9. Show if \( G \) is a group in which \((ab)^n = a^nb^n\) for all \( a, b \in G \) and \( n = i, i+1, i+2 \) for some positive integer \( i, \) then \( G \) is abelian.

10. Let \( W \subset G. \) Show that \( \langle W \rangle \) is the unique smallest subgroup of \( G \) containing \( W. \)

11. If \( G \) is a group and \( W \) is a nonempty subset of \( G, \) a subset, show that

\[
\langle W \rangle = \{ g \in G \mid \text{There exist } w_1, \ldots, w_r \in W \text{ (some } r \text{) not necessarily distinct such that } g = w_1^{e_1} \cdots w_r^{e_r}, \ e_1, \ldots, e_r \in \mathbb{Z} \},
\]

i.e., \( W \) is an “alphabet” for the “words” (elements) in \( \langle W \rangle. \) Unfortunately, spelling may not be unique in \( G \) for any “alphabet”.

12. Let \( F \) be a field (or any non-trivial ring). Show that \( \text{GL}_n(F) \) is not abelian for any \( n > 1. \)

13. Let \( p \) be a prime and \( F = \mathbb{Z}/p\mathbb{Z}. \) Show that \( F \) is a field.

[Hint: First show that \( F \) is a domain, i.e., a commutative ring satisfying: whenever \( a, b \in F \) satisfy \( ab = 0 \) in \( F \) then \( a = 0 \) or \( b = 0. \) Then show that any domain with finitely many elements is a field.]

14. Compute \((\mathbb{Z}/8\mathbb{Z})^\times). \) Find a different group isomorphic to it. Can you do the same for \((\mathbb{Z}/2^n\mathbb{Z})^\times\) if \( n = 4, 5? \) if \( n \) is any integer at least three?

15. Let \( \mathcal{F} := \{ f \mid f : \mathbb{Z}^+ \to \mathbb{C} \} \) and \( I \) in \( \mathcal{F} \) be given by \( I(1) = 1 \) and \( I(n) = 0 \) for all \( n > 1. \) Show that \( \mathcal{F} \) is an abelian monoid and
\[ \mathcal{A} := \{ f \mid f \in \mathcal{F} \text{ with } f(1) \neq 0 \} \]
is an abelian group with unity \( I \) under the Dirichlet product given by
\[
(f \ast g)(n) := \sum_{d|n} f(d)g(n/d).
\]

9. First Properties

Recall that we defined a subgroup of a group as follows:

**Definition 9.1.** A subset \( H \) of a group \( G \) is called a subgroup of \( G \) if it becomes a group under the restriction of the binary operation on \( G \) to \( H \), i.e., \( \cdot|_{H \times H} : H \times H \to H \) makes sense (meaning that the image of \( \cdot|_{H \times H} \) lies in \( H \)), and has the same unity as \( G \).

We begin this section with a criterion for a subset of a group to be a subgroup. This can be compared to the analogous criterion for sets to be subspaces of a vector space.

**Proposition 9.2.** Let \( G \) be a group and \( H \) be a non-empty subset of \( G \). Then \( H \) is a subgroup of \( G \) if and only if the following two conditions hold:

(i) If \( a \) and \( b \) are elements of \( H \), then \( ab \) is an element of \( H \).

[We say that \( H \) is closed under \( \cdot \), i.e., this means that the restriction
\[ \cdot|_{H \times H} : H \times H \to H \]
is defined. Note we have changed the target from \( G \) to \( H \).]

(ii) If \( a \in H \) then \( a^{-1} \in H \).

Moreover, these two conditions are equivalent to

\[(*) \quad \text{If } a, b \in H \text{ then } ab^{-1} \in H.\]

**Proof.** \((\Rightarrow)\) follows from the definition of subgroup, since by definition, \( e_G = e_H \).

\((\Leftarrow)\): We first note that \((i)\) and \((ii)\) imply \((*)\) for if \( a, b \in H \), then \( a, b^{-1} \in H \), hence \( ab^{-1} \in H \). Conversely, suppose that \((*)\) holds. As \( a \in H \) implies \( a, a \in H \), we have \( e_G = aa^{-1} \in H \). By definition, \( e_G a = a = ae_G \) for all \( a \in H \), so \( e_G = e_H \). As \( e_G, a \in H \), we have \( e_G a^{-1} = a^{-1} \) lies in \( H \). We conclude that if \( a, b \in H \), then \( a, b^{-1} \in H \), hence \( ab = a(b^{-1})^{-1} \in H \). Finally, since associativity holds in \( G \), it holds in the subset \( H \).

Note that property \((*)\) implies that we do not need to assume in the definition of a subgroup that it has the same unity as the group.
Corollary 9.3. Let $G$ be a group and $H \subseteq G$ a non-empty subset. If $H$ is a finite set, then $H$ is a subgroup if and only if $H$ is closed under $\cdot$.

**Proof.** We need only check:

$(\Leftarrow)$: It suffices to show if $a \in H$ then $a^{-1} \in H$. Our hypothesis implies that

$$S := \{a^n \mid n \in \mathbb{Z}^+\} \subseteq H$$

and is also a finite set. We need the following, whose proof we leave as an exercise:

**Dirichlet’s Pigeonhole Principle:** If $N \geq n + 1$ objects are placed in $\leq n$ boxes, then at least one box contains at least two elements.

In our case, $|S| < \infty$, so there exist integers $n > m \geq 1$ satisfying $a^n = a^m$. Therefore, $a^{n-m} = a^n a^{-m} = e_G = aa^{n-m-1}$ in $G$. It follows that $e_G$ lies in $H$ and $a^{n-m-1} = a^{-1}$ lies in $H$ if $n - m - 1 > 0$. \hfill \square

Note that this proof shows that if $G$ is a group with $a \in G$ and $\langle a \rangle$ finite, then there exists $N \in \mathbb{Z}^+$ such that $a^N = e$ and by well-ordering there exists a minimal such $N$.

If $G$ is a group, we call $|G|$ the order of $G$, so $G$ is a finite group if it has finite order, infinite otherwise. If $a \in G$, we call $|\langle a \rangle|$ the order of $a$ and say $a$ has finite order if $|\langle a \rangle|$ is finite.

**Examples 9.4.**

1. $(\mathbb{Z}, +) = \langle 1 \rangle$ is an infinite cyclic group.

2. If $n \in \mathbb{Z}^+$ then $(\mathbb{Z}/n\mathbb{Z}, +) = \langle 1 \rangle$ is a cyclic group of order $n$ as is $
\mu_n = \langle e^{2\pi \sqrt{-1}/n} \rangle$.

3. $\mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2\mathbb{Z}$ is an abelian group of order four and not cyclic.

4. If $n > 2$ then $D_n$ is a non-abelian group of order $2n$.

5. The quaternion group is a non-abelian group of order 8. [Does it look the same algebraically as $D_4$?]

6. If $n > 2$, then $S_n$ is a non-abelian group of order $n!$.

7. $(\mathbb{Z}/n\mathbb{Z})^\times$ is an abelian group of order $\varphi(n)$.

8. $GL_n(\mathbb{R})$, $n > 1$, is an infinite non-abelian group.

We need to know how to determine when two groups are the same algebraically.

Let $\varphi : G \to G'$ be a map of groups. We call $\varphi$ a group homomorphism if

$$\varphi(a \cdot b) = \varphi(a) \cdot \varphi(b) \text{ for all } a, b \in G$$

$$\varphi(e_G) = e_{G'}.$$
Note in the first equation \( \cdot \) on the left hand side is in \( G \) and on the right hand side is in \( G' \). We shall also see that we really do not need the second equation. We put it in because “homomorphisms” of algebraic structures should preserve those structures, and the unity is a special element of a group. That it is preserved given the first equation is fortuitous but should not \textit{a priori} be assumed.

If \( \varphi \) above is a (group) homomorphism and if, in addition, \( \varphi \) is:

1. injective, we say \( \varphi \) is a (group) \textit{monomorphism} or \textit{monic}.
2. surjective, we say \( \varphi \) is a (group) \textit{epimorphism} or \textit{epic}.
3. bijective and \( \varphi^{-1} \) is a homomorphism, we say \( \varphi \) is a (group) \textit{isomorphism}.

We let
\[
\ker \varphi := \{ a \in G \mid \varphi(a) = e_{G'} \}
\]
called the \textit{kernel} of \( \varphi \) and
\[
\im \varphi := \{ \varphi(a) \in G' \mid a \in G \}
\]
called the \textit{image} of \( \varphi \).

If \( G \) and \( G' \) are groups and there exists an isomorphism \( \varphi : G \rightarrow G' \), which we also write as \( \varphi : G \xrightarrow{\sim} G' \), we say that \( G \) and \( G' \) are \textit{isomorphic} and write \( G \cong G' \).

\textbf{Remark 9.5.} Let \( \varphi : G \rightarrow G' \) be a group homomorphism. We leave it as an exercise to show that \( \varphi \) is a monomorphism if and only if given any group homomorphisms \( \psi_1, \psi_2 : H \rightarrow G \) with compositions satisfying \( \varphi \circ \psi_1 = \varphi \circ \psi_2 \), then \( \psi_1 = \psi_2 \); and \( \varphi \) is an epimorphism if and only if given any group homomorphisms \( \theta_1, \theta_2 : G' \rightarrow H \) with compositions satisfying \( \theta_1 \circ \varphi = \theta_2 \circ \varphi \), then \( \theta_1 = \theta_2 \). (Cf. Exercise 1.12(6) and (7).)

\textbf{Remark 9.6.} You should know that a bijective linear transformation of vector spaces is an isomorphism, i.e., the inverse is automatically linear. The same is true for groups, i.e., a bijective homomorphism of groups is an isomorphism. We shall see this is also true in other algebraic cases. The proof is essentially the same as that for bijective linear transformations.

\textbf{Properties 9.7.} Let \( \varphi : G \rightarrow G' \) be a group homomorphism.

1. \( \varphi(e_G) = e_{G'} \) (i.e., this is not needed in the definition of a subgroup).
2. \( \varphi(a^{-1}) = \varphi(a)^{-1} \) for all \( a \in G \).
3. \( \ker \varphi \subset G \) and \( \im \varphi \subset G' \) are subgroups.
4. \( \varphi \) is monic if and only if \( \ker \varphi = \{e_G\} \).

5. \( \varphi \) is epic if and only if \( \text{im} \varphi = G' \).

**Proof.** (1): \( e_G \varphi(e_G) = \varphi(e_G) = \varphi(e_G e_G) = \varphi(e_G) \varphi(e_G) \).

By Cancellation, \( e_G = \varphi(e_G) \).

(2): \( e_G' = \varphi(e_G) = \varphi(a a^{-1}) = \varphi(a) \varphi(a^{-1}) \). Similarly, \( e_G' = \varphi(a^{-1}) \varphi(a) \).

(3), (5) are left as exercises.

(4): We have \( \varphi(a) = \varphi(b) \) if and only if \( e_G' = \varphi(a) \varphi(b)^{-1} = \varphi(a) \varphi(b^{-1}) = \varphi(ab^{-1}) \)

if and only if \( ab^{-1} \in \ker \varphi \). So \( \ker \varphi = \{e_G\} \) if and only if \( \varphi(a) = \varphi(b) \)
implies \( a = b \). \( \Box \)

**Examples 9.8.**

1. Let \( G \) and \( H \) be groups. Then the map \( \varphi : G \to H \)
given by \( x \mapsto e_H \) is a group homomorphism, called the *trivial homomorphism*.

2. Let \( H \) be a subgroup of \( G \). Then the inclusion map of \( H \) in \( G \) is a group homomorphism. This observation leads to a better way to define the notion of a subgroup, viz., let \( H \subset G \) be a subset with both \( H \) and \( G \) groups. Then \( H \) is a subgroup of \( G \) if and only if the inclusion map is a monomorphism.

3. Let \( F \) be a field. Then the map \( \text{det} : \text{GL}_n(F) \to F^\times \) given by \( A \mapsto \text{det} A \)
is an epimorphism with \( \ker \text{det} = \text{SL}_n(F) \). (The same is true if \( F \) is just a commutative ring.)

4. Let \( n \in \mathbb{Z}^+ \). Then the map \( \overline{\cdot} : \mathbb{Z} \to \mathbb{Z}/n\mathbb{Z} \) given by \( x \mapsto \overline{x} \) is an epimorphism of additive groups with \( \ker \overline{\cdot} = n\mathbb{Z} := \{kn \mid k \in \mathbb{Z}\} \).

5. Let \( n \in \mathbb{Z}^+ \). Then the map \( \varphi : \mathbb{Z}/n\mathbb{Z} \to \mathbb{C}^\times \) given by \( \overline{x} \mapsto e^{2\pi i x/n} \) is a well-defined (why?) map, and is a monomorphism with \( \text{im} \varphi = \mu_n \).

6. Let \( G \) be an abelian group. Then the map \( \varphi : G \to G \) given by \( x \mapsto x^{-1} \) is a group homomorphism. If \( G \) is not abelian, this map is never a group homomorphism.

Classifying objects up to isomorphism is usually impossible and difficult when possible. Indeed some of the most important theorems arise when this can be accomplished. However, in a few cases it is quite elementary. For example, two finite sets are isomorphic (i.e., bijective) if and only if they have the same number of elements, and two finite dimensional vector spaces over a field \( F \) are isomorphic if and only if
they have the same dimension. Another is the case for cyclic groups that we now show.

**Theorem 9.9. (Classification of Cyclic Groups)** Let $G = \langle a \rangle$ be a cyclic group and $\theta : \mathbb{Z} \to G$ the map given by $m \mapsto a^m$. Then $\theta$ is a group epimorphism. It is an isomorphism if and only if $G$ is infinite. If $G$ is finite, then $|G| = n$ if and only if $\ker \theta = n\mathbb{Z}$. If this is the case, then the map $\overline{\theta} : \mathbb{Z}/n\mathbb{Z} \to G$ given by $\overline{m} \mapsto a^m$ induced by $\theta$ is a group isomorphism.

**Proof.** As $\theta(i + j) = a^{i+j} = a^i a^j = \theta(i)\theta(j)$, the map is a group homomorphism. It is clearly an epimorphism. We must show that $\theta$ is an isomorphism if and only if $|G|$ is infinite and if not then $\ker \theta \neq \{e\}$ to establish the first part of the theorem. We have $\theta$ is one-to-one if and only if $a^i \neq a^j$ whenever $i \neq j$ if and only if $a^m \neq e_G$ for all $m \neq 0$ if and only if $\theta$ is an isomorphism (as it is onto). In particular, if $\theta$ is one-to-one, then $G$ is infinite.

So we may assume that there exists $N \in \mathbb{Z}^+$ with $a^N = e_G$, hence $\theta$ is not one-to-one. We show that this means that $G$ is finite. By the Well-Ordering Principle, there exists a least $n \in \mathbb{Z}^+$ such that $a^n = e_G$. We show

**Claim 9.10.** $a^i = a^j$ if and only if $i \equiv j \mod n$. In particular, $a^n = e_G$ if and only if $n \mid m$.

Note that the claim implies that $|G| = n < \infty$ and $\ker \theta = n\mathbb{Z}$. So suppose that the integers $i$ and $j$ are not equal. We may assume that $j > i$. Set $m = j - i \in \mathbb{Z}^+$ and write $m = kn + r$ with $0 \leq r < n$ and $k$ integers using the Division Algorithm. We know that

$a^i = a^j$ if and only if $a^m = a^{j-i} = e_G$

if and only if $a^r = (a^n)^k a^r = a^m = e_G$.

By the minimality of $n$, this can occur if and only if $n \mid m$ and only if $m \equiv 0 \mod n$ if and only if $i \equiv j \mod n$ as required.

We know by the First Isomorphism of Sets 7.1 that $\theta$ induces a bijection $\overline{\theta}$. [Make sure that you understand this.] It is clearly a homomorphism, so a group isomorphism.

Subgroups of a cyclic group are classified by the following, whose proof we leave as an exercise [that you should do]:

**Theorem 9.11. (Cyclic Subgroup Theorem)** Let $G = \langle a \rangle$ be a cyclic group, $H \subset G$ a subgroup. Then the following are true:
(1) \( H = \{ e \} \) or \( \langle a^m \rangle \) with \( m \geq 1 \) the least positive integer such that \( a^m \in H \). If \( |G| = n (< \infty) \) then \( m \) \( | n \). If \( G \) is infinite then \( |H| = 1 \) or \( H \) is infinite.

(2) If \( |G| = n \) and \( m \) \( n \) then \( \langle a^m \rangle \) is the unique subgroup of \( G \) of order \( \frac{n}{m} \).

(3) If \( |G| = n \) and \( m \not| n \) then \( G \) has no subgroup of order \( m \).

(4) If \( |G| = n \), the number of subgroups of \( G \) is equal to the number of positive divisors, \( \sum_{d|n} 1 \), of \( n \) (where in such a sum, we always assume \( d > 0 \)).

(5) If \( |G| \) is a prime then \( \{ e \} \) and \( G \) are the only subgroups of \( G \).

Exercises 9.12.

1. Let \( p \) be a prime and \( F = \mathbb{Z}/p\mathbb{Z} \). Then \( F \) is a field by Exercise 8.5 (13). Compute \( |G| \) for \( G = \text{GL}_n(F), \text{SL}_n(F), \text{T}_n(F), \text{ST}_n(F), \) and \( D_n(F) \).

2. Prove Dirichlet’s Pigeonhole Principle.

3. Prove Remark 9.5.

4. Prove Property 9.7(3) and Remark 9.6, i.e., a bijective group homomorphism is a group isomorphism. In particular, if \( \varphi : G \to G' \) is a monomorphism then \( G \) is isomorphic to \( \varphi(G) \).

5. Let \( G_i, i \in I \), be groups. Show that \( \times_I G_i \) satisfies the following property relative to the group homomorphisms \( \pi_j : \times_I G_i \to G_j \) defined by \( \{ g_i \}_I \mapsto g_j \) for all \( j \in I \): Given a group \( G \) and group homomorphisms \( \phi_j : G \to G_j \), there exists a unique group homomorphism \( \psi : G \to \times_I G_i \) satisfying \( \phi_j = \pi_j \circ \psi \) for all \( j \in I \).

6. Prove that every group of order four is either isomorphic to \( \mathbb{Z}/4\mathbb{Z} \) or \( \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2\mathbb{Z} \).

7. Prove the Cyclic Subgroup Theorem (Theorem 9.11).

8. Let \( G \) be an abelian group. Suppose that the elements \( a, b \) of \( G \) are of relatively prime orders \( m \) and \( n \) respectively. Show \( ab \) has order \( mn \). Is this true if \( G \) is not abelian? Prove or give a counterexample.

9. Find a group isomorphic to \( \text{ST}_2(\mathbb{R}) \) besides itself. Prove your answer.

10. Show that the quaternion group \( Q \) has a unique subgroup of order two, in particular, it is not isomorphic to \( D_4 \).

11. Let \( G \) be the subgroup of \( \text{GL}_2(\mathbb{C}) \) generated by

\[
\begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}, \begin{pmatrix} \sqrt{-1} & 0 \\ 0 & -\sqrt{-1} \end{pmatrix}, \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix}, \text{ and } \begin{pmatrix} 0 & \sqrt{-1} \\ \sqrt{-1} & 0 \end{pmatrix}.
\]

Show \( G \) is isomorphic to the quaternion group \( Q \). 

10. Cosets

We now define an equivalence on a group $G$ that mimicks congruence modulo $m$ for integers. Although the similarity is precise, the reason for doing so is not so clear. Later we shall see how this equivalence arises in a natural way.

Let $H$ be a subgroup of $G$. Define $\equiv \mod H$ by

$$a \equiv b \mod H \text{ if and only if } b^{-1}a \in H.$$ 

Then $\equiv \mod H$ is an equivalence relation. [Show this.] If $a \in G$, the equivalence class $\overline{a}$ of $a$ is called the left coset of $a$ relative to $H$. Analogous to $\equiv \mod m$ for integers, we shall usually write $aH$ for $\overline{a}$ — remember the binary operation for $\mathbb{Z}$ is $+$ — and shall write $G/H$ for $\overline{G} = G/\equiv \mod H$. So we have the natural set surjection

$$- : G \to G/H \text{ given by } a \mapsto \overline{a} = aH.$$ 

Warning 10.1. In general, $G/H$ is not a group, so the natural surjection above is not a group homomorphism. [Can you give an example?]

Observation 10.2. Let $H \subset G$ be a subgroup and $H = eH$ ($e = e_G = e_H$) and $a \in G$. Then

$$aH = H \text{ if and only if } aH = H = eH \text{ if and only if } a = e^{-1}a \in H,$$

i.e., in our other notation

$$\overline{a} = \overline{e} \text{ if and only if } e \in \overline{a} \text{ if and only if } a \in \overline{e}.$$ 

We also have

$$aH = \{b \in G \mid b \equiv a \mod H\} = \{b \in G \mid a^{-1}b \in H\}$$

$$= \{b \in G \mid a^{-1}b = h \text{ some } h \in H\}$$

$$= \{b \in G \mid b = ah \text{ some } h \in H\} = \{ah \mid h \in H\}.$$ 

(Cf. $a + m\mathbb{Z}$ for the equivalence $\equiv \mod m$.)

We now use the Mantra of Equivalence Relations 5.13. Let $H$ be a subgroup of $G$ and $\mathcal{H}$ be a system of representatives for $\equiv \mod H$. We call $\mathcal{H}$ a left transversal of $H$ in $G$. Then we have

Mantra 10.3. for Cosets. In the above setup, we have

$$G = \bigvee_{\mathcal{H}} aH.$$ 

In particular, if $|G| < \infty$, then

$$|G| = \sum_{\mathcal{H}} |aH|.$$
As mentioned before, this is only useful if we can compute the size of the equivalence class $aH$. We call $|H|$ the index of $H$ in $G$ and denote it by $[G : H]$. The Mantra allows us to establish our first important counting result:

**Theorem 10.4.** (Lagrange’s Theorem) Let $G$ be a finite group and $H$ a subgroup of $G$. Then


**Proof.** We begin with the following:

**Claim 10.5.** If $G$ is an arbitrary group (i.e., without assuming it is finite) and $H$ a subgroup, then $|aH| = |H|$ for all $a \in G$.

Define the left translation map

$$\lambda_a : H \to aH \quad by \quad h \mapsto ah.$$ 

By our computation above, we know that $\lambda_a$ is onto. But it is one-to-one also, as $ah = ah'$ in $G$ with $h, h' \in G$ implies that $h = h'$ by cancellation. Thus $\lambda_a$ is a bijection, establishing the claim.

If $\mathcal{H}$ is a left transversal for $H$ in $G$, the mantra together with the claim imply that

$$|G| = \sum_{\mathcal{H}} |aH| = \sum_{\mathcal{H}} |H| = |\mathcal{H}||H|$$

as needed. □

**Notation 10.6.** If $G$ is a group and $H = \{e_G\}$, we shall write 1 for $H$ unless $G$ is additive in which case we shall write 0 for $H$. The subgroup $\{e_G\}$ is called the trivial group in $G$.

With this notation, $|G| = |G : 1|$, so if $G$ is finite, Lagrange’s Theorem can be written

$$[G : 1] = [G : H][H : 1].$$

We, of course, have an analogous Lagrange’s Theorem for right cosets. [Definition?] It follows by both Lagrange’s Theorems (i.e., for left and right cosets) that if $G$ is a finite group and $H$ a subgroup then $|G|/|H|$ is equal to the (left) index of $H$ in $G$ which is equal to the right index of $H$ in $G$, i.e., the number of right cosets of $H$ in $G$ is the same as the number of left cosets of $H$ in $G$. This means that $[G : H]$ makes sense without prescribing right or left cosets. It does not mean, however, that if $a \in G$, then $aH = Ha$. This is usually false. We shall find an important condition on $H$ for this to be true for all $a$ in $G$. 
Warning 10.7. In general, the converse to Lagrange’s Theorem is false, i.e., if $G$ is a finite group with $m \in \mathbb{Z}^+$ satisfying $m \mid |G|$, then there may not be a subgroup $H$ in $G$ such that $|H| = m$, if $m \neq 1$ or $|G|$. Of course, if $G$ is a finite cyclic group then the converse does hold by the Cyclic Subgroup Theorem 9.11.

Lagrange’s Theorem has many immediate consequences.

Corollary 10.8. Let $G$ be a finite group and $a \in G$. Then the order of $a$ divides $|G|$.

Corollary 10.9. Let $G$ be a group and $K$ and $H$ two finite subgroups of $G$ of relatively prime order. Then $K \cap H = \{e_G\}$.

Proof. We leave this as an exercise. □

Corollary 10.10. Let $G$ be a finite group of prime order $p$. Then $G \cong \mathbb{Z}/p\mathbb{Z}$. In particular, $1$ and $G$ are the only subgroups of $G$.

Notation 10.11. If $A$ is a subset of $B$ but $A \neq B$, we shall write $A < B$.

The corollary says if $G$ is a finite group of prime order, then $G$ contains no proper subgroups, i.e., $G$ contains no subgroup $H$ satisfying $1 < H < G$.

Corollary 10.12. Let $G$ be a finite group and $a \in G$. Then $a^{|G|} = e$.

Proof. If $n$ is the order of $\langle a \rangle$, then $|G| = nm$ for some $m \in \mathbb{Z}^+$ so $e = (a^n)^m = a^{|G|}$. □

Corollary 10.13. (Euler’s Theorem) Let $m, n$ be relatively prime integers with $m > 1$. Then $n^{\varphi(m)} \equiv 1 \mod m$, where $\varphi$ is the Euler $\varphi$-function.

Proof. We have shown that the unit group $(\mathbb{Z}/m\mathbb{Z})^\times$ is given by

$$(\mathbb{Z}/m\mathbb{Z})^\times = \{\bar{a} \mid a \in \mathbb{Z}, (a,m) = 1\}$$

and by definition, its cardinality is $\varphi(m)$. □

Corollary 10.14. (Fermat’s Little Theorem) Let $p$ be a positive prime integer. Then $n^p \equiv n \mod p$ for all integers $n$. If $p \nmid n$, then $n^{p-1} \equiv 1 \mod p$.

Proof. As $\varphi(p) = p - 1$, this follows from Euler’s Theorem together with the observation that $0^p = 0$. □

We give a simple application of Euler’s Theorem that shows the usefulness of elementary number theory in coding theory. All messages
to be encoded are assumed to have been translated into (sequences of) integers. This application is called the RSA code that allows encoding and decoding using a public key, i.e., an integer \( m \) publicly known, together with a secret key, a number not publicly known. It and its variants are a major way of sending important data, e.g., credit card numbers, over the internet. To be effective, the secret key is based on the fact that the public key cannot be factored (in a reasonable amount of time). To set up the code, one chooses the public key \( m \) to be an integer that is the product of two large distinct primes \( p \) and \( q \). We know that

\[
\varphi(m) = \varphi(p)\varphi(q) = (p-1)(q-1) = pq - (p + q) + 1 = m - (p + q) + 1.
\]

So to know \( \varphi(m) \), we must know \( p \) and \( q \), hence we must be able to factor \( m \). A public encoder \( e \) is chosen and fixed. This is an integer satisfying \((e, \varphi(m)) = 1\). Then \( d \) is chosen (using the knowledge of \( m = pq \), hence \( \varphi(m) \)) to be an integer satisfying \( de \equiv 1 \mod \varphi(m) \), i.e., the inverse of \( e \) modulo \( \varphi(m) \). Write \( de = 1 + k\varphi(m) \) with \( k \in \mathbb{Z} \). \( d \) is the (non-public) decoder. Messages are written in numbers \( x < m \).

**Claim 10.15.** \( x^{ed} \equiv x \mod m \), i.e., for all positive integers \( x < m \), we get \( x \) back.

If \((x, m) = 1\), then by Euler’s Theorem, we have

\[
x^{\varphi(m)} = x \cdot x^{k\varphi(m)} \equiv x \mod m
\]

as needed, so we may assume that \((x, m) \neq 1\). As \( x \) satisfies \( 1 < x < m \), either \( p \mid x \) or \( q \mid x \) but not both. We may assume that \( p \mid x \). Therefore, \( x = pn \) for some integer \( n \) and \((x, q) = 1\). Since \( p \mid x \),

\[
x^{ed} \equiv 0 \equiv x \mod p.
\]

By Fermat’s Theorem,

\[
x^{\varphi(m)} = (x^{q-1})^{p-1} \equiv 1^{p-1} \equiv 1 \mod q.
\]

Since \( p \) and \( q \) are relatively prime and both \( p \) and \( q \) divide \( x^{ed} - x \), we have \( m = pq \mid x^{ed} - x \), so \( x^{ed} \equiv x \mod m \). This establishes the claim.

Some care must be taken in the choice of the large primes \( p \) and \( q \). For example, if all primes less than \( N \) are known and either \( p - 1 \) or \( q - 1 \) factors into products of such primes, then \( m \) can be factored using
what is called the Pollard \((p - 1)\)-Factoring Algorithm. Note that the idea of the RSA code is also based on group theory. Other public codes have been developed using group theory.

**Exercises 10.16.**

In the following exercises, let \(G\) be a group with \(H\) and \(K\) subgroups.

1. Show that \(\equiv \mod H\) is an equivalence relation.

2. Let \(G\) be a group of order \(p^n\) where \(p\) is a prime and \(n \geq 1\). Prove that there exists an element of order \(p\) in \(G\).

3. Classify all groups of order eight up to isomorphism.

4. Let \(HK := \{hk \mid h \in H, k \in K\}\). Then (clearly) \(H/(H \cap K)\) is a subset of \(G/K\). Show that \(f : H/(H \cap K) \to HK/K\) by \(h(H \cap K) \mapsto hK\) is a well-defined bijection. In particular, if \(G\) is a finite group, then \(|HK||H \cap K| = |H||K|\).

5. Suppose that \(G\) is a finite group. Show \([G : H \cap K] \leq [G : K][G : H]\) with equality if \([G : H]\) and \([G : K]\) relatively prime.

6. (Poincaré) Suppose both \(H\) and \(K\) have finite index in \(G\). Show that \(H \cap K\) has finite index in \(G\).

7. If \(K \subset H \subset G\), show that \([G : K] = [G : H][H : K]\) (even if any are infinite if read correctly).

8. If \(K\) and \(H\) are finite subgroups of \(G\) of relatively prime degree, show that \(K \cap H = \{e_G\}\).

9. Let \(a \in (\mathbb{Z}/n\mathbb{Z})^\times\) have order \(m\) and \(b \in (\mathbb{Z}/n\mathbb{Z})^\times\) have order \(n\).
   (i) Show that the order of \(a^r\) is \(r/(m, r)\).
   (ii) Show if \((m, n) = 1\), then the order of \(ab\) is \(mn\).

10. Let \(p\) be an odd prime. Prove that \((\mathbb{Z}/p\mathbb{Z})^\times\) is a cyclic group of order \(p - 1\). [You may use the fact that the polynomial \(t^n - 1\) has at most \(n\) roots in \((\mathbb{Z}/p\mathbb{Z})[t]\) for any \(n > 0\).]

11. Let \(m\) and \(n\) be positive integers that are not relatively prime. Show that \(\mathbb{Z}/m\mathbb{Z} \times \mathbb{Z}/n\mathbb{Z}\) is abelian but not cyclic.

12. Prove that a group of order 30 can have at most 7 subgroups of order 5.

13. (Wilson’s Theorem) Let \(p\) be a positive integer. Then \(p\) is a prime if and only if \((p - 1)! \equiv -1 \mod p\). (This is not a practical criterion.) [Hint: If \(1 \leq j \leq p - 1\), when can \(j^2 \equiv 1 \mod p\)?]

### 11. Homomorphisms

We begin with a useful fact whose proof is left as an exercise:
Observation 11.1. If $\varphi : G \to G'$ is a group homomorphism and $H$ is a subgroup of $G$, then $\varphi(H) \subset G'$ is a subgroup.

If $G$ is a group and $x \in G$, then the map $\theta_x : G \to G$ given by $g \mapsto xgx^{-1}$ is called conjugation by $x$. If $H \subset G$ is a subgroup, then we let $xHx^{-1} := \theta_x|_H(H) = \theta_x(H) = \{xhx^{-1} \mid h \in H\}$.

Lemma 11.2. Let $G$ be a group and $x \in G$. Then $\theta_x : G \to G$ is an isomorphism. In particular, if $H \subset G$ is a subgroup, so is $\theta_x(H)$ and $H \cong \theta_x(H) = xHx^{-1}$. In particular, $|H| = |xHx^{-1}|$. Since $\theta_x \theta_x^{-1} = 1_G = \theta_x^{-1} \theta_x$.

Proof. The map $\theta_x$ is a bijection with inverse $\theta_x^{-1}$. Let $g, g' \in G$. Then the equation

\[(11.3) \quad \theta_x(gg') = x gg' x^{-1} = x g e g' x^{-1} = x g x^{-1} x g' x^{-1} = \theta_x(g) \theta_x(g')\]

shows that $\theta_x$ is a homomorphism. $\Box$

This idea of writing an identity in a clever way is a most useful device, which we shall call Great Trick. You have undoubtedly used it quite often for the elements 1 and 0 in $\mathbb{R}$.

By Observation 11.1, we know that $\theta_x(H) \subset G$ is a subgroup, so we can view $\theta_x : H \to \theta_x(H)$. As it is a bijective homomorphism, it is an isomorphism.

A group isomorphism $\varphi : G \to G$ is called a group automorphism. So $\theta_x$ above is a group automorphism for all $x \in G$. If $G$ is an abelian group, then $\theta_x$ is the identity on $G$ for all $x \in G$. If $G$ is not abelian, then there must exist elements $x$ and $y$ in $G$ satisfying $xy \neq yx$, so $\theta_x \neq 1_G$. In general, if $H$ is a subgroup, $xHx^{-1} \neq H$, i.e., the isomorphism $\theta_x|_H : H \to xHx^{-1}$ (replacing the target with $\theta_x(H)$) is not an automorphism. [Can you come up with examples when $\theta_x|_H$ is an automorphism?] Of course, if $x \in H$ then $xHx^{-1} = \theta_x(H) = H$.

We say that a subgroup $H$ of a group $G$ is normal and write $H \triangleleft G$ if $xHx^{-1} = \theta_x(H) = H$ for all $x \in G$.

Remarks 11.4. Let $H$ be a subgroup of $G$.

1. $H \triangleleft G$ does not mean that $\theta_x|_H = 1_H : H \to H$ for all $x \in G$ (or even all $x \in H$). In fact, it usually is not. For example, (see below)

   $H := \langle r \rangle \triangleleft D_3$ but $f r f^{-1} = r^{-1}$, so $\theta_f|_H \neq 1_H$. 


2. Let

\[ \text{Aut}(G) := \{ \sigma : G \to G \mid \sigma \text{ is an automorphism} \}. \]

Then \( \text{Aut}(G) \) is a group under composition. [Check this.] In fact, \( \text{Aut}(G) \) is a subgroup of \( \Sigma(G) \). A conjugation \( \theta : G \to G \), with \( x \in G \), is also called an inner automorphism of \( G \). Set

\[ \text{Inn}(G) := \{ \theta_x \mid x \in G \}. \]

Check that this is a subgroup of \( \text{Aut}(G) \), called the inner automorphism group of \( G \). Moreover, the following are equivalent:

(a) \( H \triangleleft G \).
(b) \( \theta(H) = H \) for all \( \theta \in \text{Inn}(G) \).
(c) \( \theta|_H \in \text{Aut}(H) \) for all \( \theta \in \text{Inn}(G) \).
(d) The restriction map \( \theta|_H : \text{Inn}(G) \to \text{Aut}(H) \) is defined (i.e., the image lies in \( \text{Aut}(H) \)).

3. \( H \triangleleft G \) if and only if \( xH = Hx \) for all \( x \in G \). (Why?)

4. If \( xHx^{-1} \subset H \) for all \( x \in G \), then \( y^{-1}Hy \subset H \) for all \( y \in G \) (let \( y = x^{-1} \)), so \( H \subset yHy^{-1} \) for all \( y \in G \). Hence

\[ H \triangleleft G \] if and only if \( xHx^{-1} \subset H \) for all \( x \in G \).

Examples 11.5. Let \( G \) be a group.

1. \( 1 \triangleleft G \) and \( G \triangleleft G \). If \( G \) is a non-trivial group and these are the only normal subgroups of \( G \), then we say \( G \) is a simple group. We have seen that \( \mathbb{Z}/p\mathbb{Z} \) is a simple group for all primes \( p \). One of the great theorems in mathematics is the classification of all finite simple groups. The proof is thousands of pages long.

2. If \( G \) is abelian, every subgroup is normal. The converse is false, i.e, there exist nonabelian groups in which every subgroup is normal.

3. Let

\[ Z(G) := \{ x \in G \mid xg = gx \text{ for all } g \in G \}. \]

This set is called the center of \( G \). It is a normal subgroup of \( G \). In fact, any subgroup of \( Z(G) \) is a normal subgroup of \( G \). Moreover, \( G \) is abelian if and only if \( G = Z(G) \).

4. If \( F \) is a field (even a commutative ring), then \( \text{SL}_n(F) \triangleleft \text{GL}_n(F) \).

5. Let \( \theta : G \to G' \) be a group homomorphism. Then \( \ker \theta \triangleleft G \).

[In general, \( \text{im} \theta \) is not a normal subgroup of \( G \). [Can you give an example?]
6. If $G$ is an abelian group, then $\varphi : G \to G$ given by $x \mapsto x^{-1}$ is an automorphism. It is the identity map if and only if $x^2 = e$ for all $x \in G$, otherwise it is of order two. In particular, if $G$ is a finite abelian group of order at least three, then $2 \mid |\text{Aut}(G)|$. [Can you prove this?]

7. $\text{Inn}(G) \triangleleft \text{Aut}(G)$.

8. If $H$ is a subgroup of $G$ of index two then $H \triangleleft G$. Indeed, if $a \in G \setminus H$ then $G = H \lor aH = H \lor Ha$, so we must have $aH = Ha$, hence $aHa^{-1} = H$.

For the next example, we shall need matrix representations of linear transformations of finite dimensional vector spaces. In Appendix B, we review this notation. (This appendix does a more general case but is applicable here.) In particular, if $T : V \to W$ is a linear transformation of vector spaces over a field $F$ with bases $\mathcal{B}$ and $\mathcal{C}$ respectively, then we let $[T]_{\mathcal{B},\mathcal{C}}$ denote the matrix representation of this linear transformation relative to these bases. If $V = V'$ and $\mathcal{B} = \mathcal{C}$, we write this as $[T]_{\mathcal{B}}$.

9. If $\sigma \in S_n$, we can write $\sigma$ as

$$
\begin{pmatrix}
1 & 2 & \ldots & n \\
\sigma(1) & \sigma(2) & \ldots & \sigma(n)
\end{pmatrix}
$$

where the top row is the elements of the domain and the bottom row the corresponding values. Let $S_n := \{e_1, \ldots, e_n\}$ be the standard basis for $V = \mathbb{R}^n$. For each $\sigma$, define a linear transformation $T_\sigma : V \to V$ by $\sum \alpha_i e_i \mapsto \sum \alpha_i e_{\sigma(i)}$. Then $T_\sigma$ is a vector space isomorphism with inverse $T_\sigma^{-1}$. Define $\theta$ by

$$
(11.6) \quad \theta : S_n \to \text{GL}_n(\mathbb{R}) \text{ is given by } \sigma \mapsto [T_\sigma]_{S_n}.
$$

It is easy to see that $\theta$ is a group monomorphism. Each $[T_\sigma]_{S_n}$ is a permutation matrix, i.e., a matrix having exactly one entry of 1 in each row and each column and 0s elsewhere. So it is just a permutation of the rows (or columns) of the identity matrix. Let $\text{Perm}_n(\mathbb{R})$ be the set of permutation matrices. It is the image of $\theta$, so a group. Changing the target, we view $\theta : S_n \to \text{Perm}_n(\mathbb{R})$, an isomorphism. If $A \in \text{Perm}_n(\mathbb{R})$, then $\det A = \pm 1$, so we have a composition of maps

$$
S_n \xrightarrow{\theta} \text{Perm}_n(\mathbb{R}) \xrightarrow{\det} \{\pm 1\}.
$$

It is a homomorphism as both $\theta$ and $\det$ are. If $n > 1$ then the kernel of this map is called the alternating group on $n$ letters and denoted by $A_n$. Elements of $A_n$ are called even permutations and
elements of $S_n \setminus A_n$ are called *odd permutations*. We have $A_n \triangleleft S_n$, and, in fact, $[S_n : A_n] = 2$ if $n \geq 2$. We shall study this group in Section 24.

**Remarks 11.7.** 1. We have seen that a group of prime order is a simple group, i.e., a group without any proper normal subgroups. These turn out to be the only abelian simple groups. The group $A_5$ is the non-abelian simple group of smallest order. This was shown by Abel. In fact,

**Theorem 11.8.** (Abel’s Theorem) *For every* $n \geq 5$, the group $A_n$ *is simple.*

The group $A_2$ is the trivial group and $A_3$ is the cyclic group of order three so simple. The group $A_4$ is not simple.

Abel’s Theorem is important as it is the key to proving the Abel-Ruffini Theorem, which states that there is no formula for the roots of a general fifth degree polynomial with rational coefficients involving only the extraction of $n$th roots and $+$ and $\cdot$. We shall prove Abel’s Theorem in Theorem 24.13.

2. If $K$ and $H$ are subgroups of $G$ satisfying $K \subset H \subset G$ with $K \triangleleft H$ and $H \triangleleft G$, it is not necessarily true that $K \triangleleft G$, i.e., being normal is not transitive.

3. If $K$ and $H$ are subgroups of $G$ satisfying $K \subset H \subset G$ with $K \triangleleft G$, then it is true that $K \triangleleft H$.

4. There is a useful stronger property that a subgroup of a group can satisfy than being normal. We say that a subgroup $H$ of a group $G$ is a characteristic subgroup of $G$ if for every $\sigma \in \text{Aut}(G)$, we have $\sigma|_{H}$ lies in $\text{Aut}(H)$. Equivalently, the restriction map $\text{res} : \text{Aut}(G) \to \text{Aut}(H)$ given by $\theta \mapsto \theta|_{H}$ is well-defined, i.e., the target is correct. If $H$ is a characteristic subgroup of $G$, we write $H \triangleleft\triangleleft G$. Clearly, unlike being normal, being characteristic is transitive, i.e., if $K \triangleleft\triangleleft H \triangleleft\triangleleft G$, then $K \triangleleft\triangleleft G$.

**Exercises 11.9.**

2. Find all subgroups of $S_3$ and determine which ones are normal.
3. Show that a subgroup $H \subset G$ is normal if and only if $gH = Hg$ for all $g \in G$. If $H$ is not normal is it still true that for each $g \in G$ there is an $a \in G$ such that $gH = Ha$?
4. Determine all subgroups of the quaternion group of order 8 and determine which ones are normal.
5. Classify all groups of order eight up to isomorphism.

6. Let $G$ be a group. Show
   (i) The center $Z(G)$ is a subgroup $G$.
   (ii) Any subgroup of the center of $G$ is a normal subgroup of $G$.
   (iii) $G$ is abelian if and only if it equals its center.

7. Determine the center of the dihedral group $D_n$ and prove your assertion.

8. Let $H$ be a subgroup of $G$. In Exercises 8.5(7) and (8), we defined the core of $H$ in $G$ to be $\text{Cor}_G(H) := \bigcap xHx^{-1}$ and the normalizer of $H$ in $G$ to be $\text{N}_G(H) := \{ x \in G \mid xHx^{-1} = H \}$, respectively. Show
   (i) $\text{Cor}_G(H)$ is the unique largest normal subgroup of $G$ contained in $H$.
   (ii) $\text{N}_G(H)$ is the unique largest subgroup of $G$ containing $H$ as a normal subgroup.

9. Show $\text{Inn}(G) \triangleleft \text{Aut}(G)$.

10. Prove the assertions in Example 11.5(6).

11. Let $G$ be a cyclic group. Determine $\text{Aut}(G)$ and $\text{Inn}(G)$ up to isomorphism as groups that we know. Prove your result.
   [Hint. Where do generators go?]

12. Let $G$ and $H$ be finite cyclic groups of order $m$ and $n$, respectively. Show the following:
   (i) If $m$ and $n$ are relatively prime, then $\text{Aut}(G \times H) \cong \text{Aut}(G) \times \text{Aut}(H)$ and is abelian.
   (ii) If $m$ and $n$ are not relatively prime, then $\text{Aut}(G \times H)$ is never abelian. (Cf. Exercise 11.9(13).)

13. Compute $\text{Aut}((\mathbb{Z}/p\mathbb{Z})^n)$ up to isomorphism and its order when $p$ is a prime.

14. Let $\varphi : G \to \text{Aut}(H)$, write $x \mapsto \varphi_x$, be a group homomorphism where $G$ and $H$ are groups. Define $H \rtimes_\varphi G$ to be the cartesian product $H \times G$ with group structure induced by the binary operation
   \[(h, g) \cdot (h', g') = (h\varphi_g(h'), gg')\]
   for all $g, g' \in G$ and $h, h' \in H$. Show that $H \rtimes_\varphi G$ is a group, called the \textit{(external) semidirect product} of $H$ and $G$, with $H \times 1$ a normal subgroup. Note that the product $G \times H$ is the case that $\varphi$ is the trivial homomorphism.

15. Let $\sigma \in \text{Aut}(\mathbb{Z}/n\mathbb{Z})$ be the automorphism defined by $x \mapsto -x$ and $\varphi : \mathbb{Z}/2\mathbb{Z} \to \text{Aut}(\mathbb{Z}/n\mathbb{Z})$ be the group homomorphism defined by $1 \mapsto \sigma$. Show that $\mathbb{Z}/n\mathbb{Z} \rtimes_\varphi \mathbb{Z}/2\mathbb{Z} \cong D_n$. 

16. Let \( p > 1 \) be a prime and \( n > 1 \). Show that \( \mathbb{Z}/p^n\mathbb{Z} \) is not isomorphic to a semidirect product.

17. Find all subgroups of \( A_4 \) and determine which ones are normal. In particular, show that \( A_4 \) is not simple.

18. Let \( G \) be a group in which every subgroup is normal. Is \( G \) an abelian group? Prove or provide a counterexample.

19. A commutator in \( G \) is an element of the form \( xyx^{-1}y^{-1} \) where \( x, y \in G \). Let \( G' \) be the subgroup of \( G \) generated by all commutators, i.e., every element of \( G' \) is the product of commutators. We call \( G' \) the commutator or derived subgroup of \( G \). It is also denoted \([G, G]\).

Show
(i) \( G' \triangleleft G \).
(ii) \( G' \triangleleft \triangleleft G \).

20. Let \( K \subset H \subset G \) be subgroups of \( G \). Show all of the following:
(i) If \( K \triangleleft \triangleleft H \) and \( H \triangleleft G \) then \( K \triangleleft G \).
(ii) \( Z(G) \triangleleft \triangleleft G \).
(iii) \( G' \triangleleft \triangleleft G \).
(iv) Inductively define \( G^{(n)} \) as follows: \( G^{(1)} = G' \). Having defined \( G^{(n)} \) define \( G^{(n+1)} := (G^{(n)})' \). Then \( G^{(n+1)} \triangleleft \triangleleft G \).

12. The First Isomorphism Theorem

Given an equivalence relation, we saw that we could always form a corresponding set of equivalence classes. If \( G \) is a group and \( H \) is a subgroup, we can form the set of equivalence classes \( G/H \). Unfortunately, in general this does not have a group structure. We also know that equivalence relations arise from surjective maps. We begin by looking at the case when \( \varphi : G \to G' \) is a group epimorphism. This defines an equivalence relation \( \sim \) on \( G \) by \( a \sim a' \) if \( \varphi(a) = \varphi(a') \) and with our usual notation, we have
\[
\overline{a} = \overline{a'} \quad \text{if and only if} \quad a \sim a' \quad \text{if and only if} \quad \varphi(a) = \varphi(a')
\]
with
\[
\overline{a} = \varphi^{-1}(\varphi(a)) = \{ x \in G \mid \varphi(x) = \varphi(a) \}.
\]
But now, we have the additional information that \( \varphi \) is a homomorphism, so \( \varphi(aa') = \varphi(a)\varphi(a') \). This means that
\[
\varphi(a) = \varphi(a') \quad \text{if and only if} \quad e_{G'} = \varphi(a)^{-1}\varphi(a') = \varphi(a^{-1}a') \quad \text{and} \quad e_{G'} = \varphi(a')\varphi(a)^{-1} = \varphi(a'a^{-1})
\]
if and only if both \( a^{-1}a' \) and \( a'a^{-1} \) lie in \( \ker \varphi \).
Let $K = \ker \varphi$. So we have

\[ \overline{a} = \overline{a'} \text{ if and only if } \varphi(a) = \varphi(a') \text{ if and only if } a^{-1}a', a'a^{-1} \in K. \]

We know that if $a^{-1}a' \in K$, then $\overline{a} = aK$, so

\[ \overline{a} = \varphi^{-1}(\varphi(a)) = \{a' \mid a' = ak \text{ some } k \in K\} = aK, \]

the left coset of $a$ in $G$ relative to $K$. Hence $\overline{a} = \varphi^{-1}(\varphi(a)) = aK$. Similarly,

\[ \overline{a} = \varphi^{-1}(\varphi(a)) = \{a' \mid a' = ka \text{ some } k \in K\} = Ka, \]

the right coset of $a$ in $G$ relative to $K$. So we now have

\[ \overline{a} = \overline{a'} \text{ if and only if } \varphi(a) = \varphi(a') \text{ if and only if } aK = a'K = \overline{a} = Ka = Ka'. \]

In particular, we must have $aK = Ka$, or equivalently, $K = aKa^{-1}$ for all $a \in G$ as $a \sim a$, i.e., $K \triangleleft G$. Of course, we knew this before, but it now indicates where cosets are coming from, at least in a special case, and ties up various mathematical threads that we are developing. Note that $e_G = \varphi^{-1}(\varphi(e_G)) = \varphi^{-1}(e_G) = \ker \varphi$.

We now apply the First Isomorphism of Sets 7.1 to obtain a commutative diagram

\[
\begin{array}{ccc}
G & \xrightarrow{\varphi} & G' \\
\downarrow & & \downarrow \\
\overline{G} & \xrightarrow{\overline{\varphi}} & \overline{G}'
\end{array}
\]

where $\overline{\varphi} : \overline{G} \to G'$ is a well-defined bijection given by $\overline{a} \mapsto \varphi(a)$. We also know that

\[ \overline{a} = \overline{a'} \text{ if and only if } a \sim a' \text{ if and only if } \overline{\varphi}(\overline{a}) = \overline{\varphi}(\overline{a'}), \]

so the fact that $\varphi$ is a homomorphism implies that

\[ \overline{\varphi}(aa') = \varphi(a)a' = \varphi(a)\varphi(a') = \overline{\varphi}(\overline{a})\overline{\varphi}(\overline{a'}). \]

We can use this and the fact that $\overline{\varphi}$ is a bijection to put a group structure on $\overline{G}$. Define

\[ \cdot : \overline{G} \times \overline{G} \to \overline{G} \text{ by } \overline{a} \cdot \overline{a'} := \overline{aa'}. \]

It is well-defined, for if $\overline{a} = \overline{\pi}$ and $\overline{a'} = \overline{\pi'}$, then

\[ \overline{\varphi}(\overline{xx'}) = \overline{\varphi}(\overline{x}) \cdot \overline{\varphi}(\overline{x'}) = \overline{\varphi}(\overline{a}) \cdot \overline{\varphi}(\overline{a'}) = \overline{\varphi}(\overline{aa'}) \]

by (*), so $\overline{xx'} = \overline{aa'}$ as $\overline{\varphi}$ is a bijection. For all $a \in G$, we have $\overline{e_G} \cdot \overline{a} = \overline{a} = \overline{a} \cdot \overline{e_G}$ and $\overline{a} \cdot \overline{a^{-1}} = \overline{e_G} = \overline{a^{-1}} \cdot \overline{a}$, so $\overline{G}$ is a group with
identity \( \ker \varphi = \overline{e_G} \) and inverse of \( \overline{a} \) given by \( \overline{a}^{-1} = a^{-1} \). Furthermore, (*) implies that
\[
\varphi(\overline{a} \cdot \overline{a'}) = \varphi(\overline{aa'}) = \varphi(\overline{a}) \cdot \varphi(\overline{a'}),
\]
so \( \varphi \) is a bijective group homomorphism, hence an isomorphism. Our diagram (†) becomes

\[
\begin{array}{ccc}
G & \xrightarrow{\varphi} & G' \\
\downarrow & & \downarrow \\
G/\ker \varphi & \xrightarrow{\overline{\varphi}} & \text{im } \varphi
\end{array}
\]

with \( G/\ker \varphi \) a group, \( \sim \) a group epimorphism and \( \overline{\varphi} \) an isomorphism.

We can generalize this to any group homomorphism \( \psi : G \rightarrow G'' \), by letting \( G' = \text{im } \psi \) and \( \varphi \) the map \( \varphi : G \rightarrow G' \) given by \( \varphi(a) = \psi(a) \). As the inclusion map \( \text{inc} : G' \rightarrow G'' \) is a group homomorphism, we are in a similar situation as in Diagram 7.3, but with all maps homomorphisms. We write this, after changing notation, as:

**Theorem 12.1.** (First Isomorphism Theorem) Let \( \varphi : G \rightarrow G' \) be a group homomorphism. Then we have a commutative diagram of groups and group homomorphisms

\[
\begin{array}{ccc}
G & \xrightarrow{\varphi} & G' \\
\downarrow & & \downarrow \text{inc} \\
G/\ker \varphi & \xrightarrow{\overline{\psi}} & \text{im } \varphi
\end{array}
\]

with \( \sim \) a group epimorphism, \( \overline{\varphi} \) a group isomorphism between \( G/\ker \varphi \) and \( \text{im } \varphi \) and \( \text{inc} \) a group monomorphism.

**Example 12.2.** Let \( G \) be a group. The map \( \theta : G \rightarrow \text{Aut}(G) \) given by \( x \mapsto (\theta_x : g \mapsto xgx^{-1}) \) is a group homomorphism. If \( \theta_x = 1_G \), the identity map on \( G \), then \( xgx^{-1} = g \), i.e., \( xg = gx \), for all \( g \in G \). Therefore, \( \ker \theta = Z(G) \), the center of \( G \). By the First Isomorphism Theorem, \( \theta \) induces an isomorphism \( \overline{\theta} : G/Z(G) \rightarrow \text{Inn}(G) \).

If \( \varphi : G \rightarrow G' \) is a group homomorphism, the First Isomorphism Theorem says that \( G/H \) is a group with \( H = \ker \varphi \). It remains to ask: If \( H \) is a subgroup of \( G \), when is \( G/H \) a group such that the canonical surjection \( \sim : G \rightarrow G/H \) a group epimorphism? This is answered by the following:

**Theorem 12.3.** Let \( G \) be a group and \( H \) a normal subgroup of \( G \). Then \( G/H \) is a group under the binary operation \( \cdot : G/H \times G/H \rightarrow G/H \).
given by \((aH, bH) \mapsto abH\). If this is the case then \(- : G \to G/H\) is a group epimorphism with kernel \(H\).

We call \(G/H\) the quotient or factor group of \(G\) by (the normal subgroup) \(H\).

**Proof.** We first show that the map \(\cdot\) is well-defined. Suppose that \(aH = a'H\) and \(bH = b'H\). We must show that \(abH = a'b'H\). Equivalently, we must show that

\[
\text{if } a^{-1}a, b^{-1}b \in H \text{ then } x = (a'b')^{-1}ab \in H.
\]

By using the Great Trick and the definition of normality, we have

\[
x = b^{-1}a^{-1}ab = b^{-1}a^{-1}a(b'b^{-1})b = (b^{-1}a^{-1}ab')(b^{-1}b)
\]

lies in \(H\) as needed.

Let \(a, b, c \in G\). Associativity holds because \(aH \cdot (bH \cdot cH) = aH \cdot bcH = abcH = abH \cdot cH = (aH \cdot bH) \cdot cH\). The unity is easily seen to be \(e_{G/H} = H = e_G H\) and \((aH)^{-1} = a^{-1}H\) for all \(a\) in \(G\). Thus \(G/H\) is a group.

The theorem shows that normal subgroups are equivalent to the kernels of group homomorphisms. Of course, it does not explain how cosets arise in a natural way if the subgroup is not a normal subgroup. We will come back to this later.

We are now in a position where we can prove some important results. Especially significant is the following:

**Theorem 12.4.** (General Cayley Theorem) Let \(H\) be an arbitrary subgroup of a group \(G\). Let \(S = G/H\), the set of left cosets of \(H\) in \(G\). For each \(x \in G\), let \(\lambda_x : S \to S\) be defined by \(gH \mapsto xgH\). Then \(\lambda_x\) is a permutation and the map \(\lambda : G \to \Sigma(S)\) given by \(x \mapsto \lambda_x\) is a group homomorphism satisfying the following two properties:

1. \(\ker \lambda \subset H\).
2. \(\ker \lambda\) is the unique largest normal subgroup of \(G\) contained in \(H\), i.e.,

\[
\text{If } N \triangleleft G \text{ and } N \subset H \text{ then } N \subset \ker \lambda.
\]

**Proof.** We first must show that the map \(\lambda_x : S \to S\), called left multiplication by \(x\), is a permutation. But this is clear since \(\lambda_{x^{-1}}\) is easily checked to be its inverse.

For all \(x, y \in G\), we have \(\lambda_{xy}(gH) = xygH = \lambda_x(\lambda_y(gH)) = \lambda_x\lambda_y(gH)\) for all \(g \in G\), so \(\lambda\) is a group homomorphism and \(\ker \lambda \triangleleft G\). If \(x \in \ker \lambda\), then \(xgH = \lambda_x(gH) = gH\) for all \(g \in G\). In particular, \(xH = H\) so \(x \in H\). Finally, suppose that \(N \triangleleft G\) with \(N \subset H\). To show \(N \subset \ker \lambda\), we must show if \(x \in N \subset H\), then \(\lambda_x\) is \(1_S\), the identity map on \(S\).
So let $x \in N$ and $g \in G$. As $g^{-1}xg \in N$, we have $xg \in gN \subset gH$. Therefore, $xgH = gH$, i.e., $\lambda_x(gH) = gH$ for all $g \in G$. Consequently, $\lambda_x = 1_S$ as required.

Therefore, $xgH = gH$, i.e., $\lambda_x(xgH) = gH$ for all $g \in G$. Consequently, $\lambda_x = 1_S$ as required.

In the above, the special case where $H = 1$ yields:

**Corollary 12.5.** (Cayley’s Theorem) Let $G$ be a group. Then the map $\lambda : G \to \Sigma(G)$ given by $x \mapsto (\lambda_x : g \mapsto xg)$ is a group monomorphism. In particular, if $G$ is a finite group of order $n$ then there exists a monomorphism $G \to S_n$.

If $S$ is a non-empty set, then a subgroup of $\Sigma(S)$ is called a permutation group. Cayley’s Theorem shows that every abstract group is isomorphic to a permutation group. The above monomorphism is called the (left) regular representation of $G$. It shows that every abstract group can be realized by a “concrete” group, e.g., a permutation group.

**Corollary 12.6.** Let $H$ be a subgroup of a group $G$ with $H < G$. If there exists no normal subgroup $N$ of $G$ satisfying $1 < N \subset H$ then $\lambda : G \to \Sigma(G/H)$ defined by $x \mapsto (\lambda_x : aH \mapsto xaH)$ is a monomorphism.

**Proof.** $\ker \lambda$ is the maximal such normal subgroup.

**Corollary 12.7.** (Useful Counting Result) Let $G$ be a finite group, $1 < H$ a subgroup of $G$ satisfying $|G| \not| |G:H|!$. Then there exists a normal subgroup $N$ of $G$ satisfying $1 < N \subset H$. In particular, $G$ is not a simple group.

**Proof.** Exercise.

We now give two applications of the Useful Counting Result. To do so we must assume the following is true (we shall prove it later).

**Theorem 12.8.** (First Sylow Theorem) Let $G$ be a finite group, $p$ a prime such that $p^s || |G|$ (i.e., $p^s | |G|$, but $p^{s+1} \not| |G|$). Then $G$ contains a subgroup of order $p^s$.

**Examples 12.9.** Let $G$ be a group.

1. If $|G| = 24$ then $G$ is not simple, as $G$ contains a group of order 8 by the First Sylow Theorem and $|G| \not| 3!$.
2. Let $p < q$ be positive primes such that $|G| = pq$. Then $G$ contains a normal group of order $q$, so is not simple.

Another useful corollary of the General Cayley Theorem is the following, whose proof we also leave as an exercise:
Corollary 12.10. Let $G$ be a finite group. Suppose that $p$ is the smallest positive prime that divides the order of $G$. Suppose also that there exists a subgroup $H$ of $G$ of index $p$. Then $H 	riangleleft G$.

Exercises 12.11.

1. Let $T : V \to W$ be a surjective linear transformation of vector spaces over $F$. Do an analogous analysis before The First Isomorphism Theorem 12.1 for $T : V \to W$ and state and prove the analogues of The First Isomorphism Theorem 12.1 and Theorem 12.3 for it.

2. Let $G$ be a finite group and $n > 1$ an integer such that $(xy)^n = x^ny^n$ for all $x, y \in G$. Let
$$G_n := \{ z \in G \mid z^n = e \} \quad \text{and} \quad G^n := \{ x^n \mid x \in G \}.$$
Show that both $G_n$ and $G^n$ are normal subgroups of $G$ and satisfy $|G^n| = [G : G_n]$.

3. Let $G$ be a finite group of order $n$, $F$ a field (or even any ring with $1 \neq 0$). Show there exists a monomorphism $\varphi : G \to \text{GL}_n(F)$, i.e., any finite group can be realized as a subgroup of matrices.

4. Prove the Useful Counting Result 12.7.

5. Prove Corollary 12.10.

6. Suppose that $G$ is a group containing a subgroup of finite index greater than one. Show that $G$ contains a normal subgroup of finite index greater than one.

7. Let $H$ be a subgroup of $G$. Define the centralizer of $H$ in $G$ to be
$$Z_G(H) := \{ x \in G \mid xh = hx \text{ for all } h \in H \}.$$
Show that it is a normal subgroup of $N_G(H)$ and the map given by $x \mapsto (\theta_x : g \mapsto xgx^{-1})$ induces $\hat{\theta} : N_G(H)/Z_G(H) \to \text{Aut}(H)$, a monomorphism defined by $xZ_G(H) \to \theta_x|_H$.

13. The Correspondence Principle

Recall if $f : A \to B$ is a set map and $D$ a subset of $B$, then the preimage of $D$ in $A$ is the set $f^{-1}(D) := \{ a \in A \mid f(a) \in D \}$. We shall need the following, whose proof is left as an exercise.

Properties 13.1. Let $f : A \to B$ be a set map, $C \subset A$, and $D \subset B$ subsets. Then:

1. $C \subset f^{-1}(f(C))$ with equality if $f$ is one-to-one.
2. $f(f^{-1}(D)) \subset D$ with equality if $f$ is onto.

We apply the First Isomorphism Theorem to establish the following important result:
**Theorem 13.2.** (Correspondence Principle) Let \( \varphi : G \to G' \) be a group epimorphism. Then

1. If \( A \) is a subgroup of \( G \) (respectively, a normal subgroup), then \( \varphi(A) \) is a subgroup of \( G' \) (respectively, a normal subgroup). In particular, group epimorphisms preserve normality.
2. If \( A \) is a subgroup of \( G \) containing \( \ker \varphi \), then \( \varphi^{-1}(\varphi(A)) = A \).
3. If \( B \) is a subgroup of \( G' \) (respectively, a normal subgroup), then \( \varphi^{-1}(B) \) is a subgroup of \( G \) (respectively, a normal subgroup) containing \( \ker \varphi \) and satisfies \( B = \varphi(\varphi^{-1}(B)) \).

In particular, if

\[ G_{\ker \varphi} := \{ A \mid A \subset G \text{ a subgroup with } \ker \varphi \subset A \} \]
\[ G' := \{ B \mid B \subset G' \text{ a subgroup} \}, \]

then \( G_{\ker \varphi} \to G' \) given by \( A \mapsto \varphi(A) \) is a bijection of sets with inverse \( B \mapsto \varphi^{-1}(B) \) preserving inclusions and restricting to a bijection on normal subgroups.

**Proof.** Let \( K = \ker \varphi \).

1. Let \( A \subset G \) be a subgroup. Then \( \varphi|_A : A \to G' \) is a group homomorphism (why?), so \( \varphi(A) = \text{im } \varphi|_A \subset G' \) is a subgroup. Next suppose that \( A \triangleleft G \) and \( y \in G' \). As \( \varphi \) is surjective, \( y = \varphi(x) \) for some \( x \in G \). Hence \( y\varphi(x)y^{-1} = \varphi(x)\varphi(A)\varphi(x)^{-1} = \varphi(xAx^{-1}) = \varphi(A) \). Therefore, \( \varphi(A) \triangleleft G' \).

2. By the properties of preimages, \( A \subset \varphi^{-1}(\varphi(A)) \), so it suffices to show if \( K \subset A \), then \( \varphi^{-1}(\varphi(A)) \subset A \). Suppose that \( x \in \varphi^{-1}(\varphi(A)) \). By definition, \( \varphi(x) = \varphi(a) \) for some \( a \in A \), hence \( e_{G'} = \varphi(a)^{-1}\varphi(x) = \varphi(a^{-1}x) \). Thus \( a^{-1}x \in K \), and so \( x \in aK \subset A \) as needed.

3. Let \( B \) be a subgroup of \( G' \). If \( x, y \in \varphi^{-1}(B) \), we have \( \varphi(xy^{-1}) = \varphi(x)\varphi(y)^{-1} \in B \), so \( \varphi^{-1}(B) \) is a subgroup of \( G \). If \( k \in K \) then \( \varphi(k) = e_{G'} \in B \), so \( K \subset \varphi^{-1}(B) \). As \( \varphi \) is onto, \( B = \varphi(\varphi^{-1}(B)) \) by the properties of preimages. Finally, suppose that \( B \triangleleft G' \). Let \( x \in A \). Then \( \varphi(x \varphi^{-1}(B)x^{-1}) = \varphi(x)\varphi(\varphi^{-1}(B))\varphi(x)^{-1} = \varphi(x)B\varphi(x)^{-1} = B \). Hence \( x\varphi^{-1}(B)x^{-1} \subset \varphi^{-1}(B) \), so \( \varphi^{-1}(B) \triangleleft G \).

We leave the following alternate form of the Correspondence Principle as an exercise. (Note we did not include the analogue of the last statement in the above.)
Theorem 13.3. (Correspondence Principle, Alternate Form) Let $G$ be a group and $K$ a normal subgroup. Let $\varphi: G \to G/K$ be the canonical epimorphism given by $x \mapsto xK$ and $L$ a subgroup of $G/K$. Then

1. There exists a subgroup $H$ of $G$ containing $K$ with $L = H/K$.
2. Let $H$ be as in (1). Then $H \trianglelefteq G$ if and only if $L \trianglelefteq G/K$.
3. Suppose that $H_1, H_2$ are two subgroups of $G$ containing $K$. If $H_1/K = H_2/K$ then $H_1 = H_2$.
4. If $G$ is a finite group and $H$ is as in (1), then $[G : H] = [G/K : H/K] = [G/K : L]$ and $|H| = |K| \cdot |L|$.

Theorem 13.4. (Third Isomorphism Theorem) Let $G$ be a group with normal subgroups $K$ and $H$ satisfying $K \subset H$. Then the map $\varphi : G/K \to G/H$ defined by $xK \mapsto xH$ is a group epimorphism with kernel $H/K$ and induces an isomorphism $\bar{\varphi} : (G/K)/(H/K) \to G/H$.

Proof. As $K$ and $H$ are normal subgroups of $G$, we know that $G/H$ and $G/K$ are groups. If $xK = yK$ then $y^{-1}x \in K \subset H$, hence $xH = yH$. Therefore, $\varphi$ is well-defined and clearly surjective. As $\varphi(xKyK) = \varphi(xyK) = xyH = xHyH = \varphi(xK)\varphi(yK)$, the map $\varphi$ is a group homomorphism. If $xK \in \ker \varphi$ then $H = e_{G/H} = \varphi(xK) = xH$, so $x \in H$, i.e., $xK \in H/K$. Conversely, if $x \in H$ then $\varphi(xK) = xH = H = e_{G/H}$, so $xK \in \ker \varphi$. The result now follows by the First Isomorphism Theorem. \qed

Let $G$ be a group and $H_1, H_2$ subgroups of $G$. Set

\[ H_1H_2 := \{h_1h_2 \mid h_i \in H_i, \ i = 1, 2\} \]
\[ = \{h_1H_2 \mid h_1 \in H_1\} = \{H_1h_2 \mid h_2 \in H_2\}. \]

By Exercise 10.16(4), we have a set bijection

(*) $f : H_1/(H_1 \cap H_2) \to H_1H_2/H_2$ given by $h_1/(H_1 \cap H_2) \mapsto h_1H_2$

for all $h_1 \in H_1$ where $H_1/(H_1 \cap H_2)$ is a subset of $G/(H_1 \cap H_2)$ and $H_1H_2/H_2$ is a subset of $G/H_2$. In particular, we have

(13.5) $|H_1H_2| = |H_1||H_2|/(H_1 \cap H_2|$

when $H_1$ and $H_2$ are finite subgroups of $G$, which will be useful in computation. We wish to determine a condition when $H_1H_2$ is a group and the set bijection $f$ in (*) is a group isomorphism. This is the content of the next isomorphism theorem. It is very useful, but harder to visualize.
**Theorem 13.6.** (Second Isomorphism Theorem) Let $G$ be a group and $H$ and $N$ be subgroups with $N$ normal in $G$. Then

1. $H \cap N \trianglelefteq H$.
2. $HN = NH$ is a subgroup of $G$.
3. $N \trianglelefteq HN$.
4. $H/H \cap N \cong HN/N$.

We have the following picture

![Diagram](https://example.com/diagram.png)

where a group below another group connected by a line is a subgroup and the quotients of the groups defined by the double lines are isomorphic.

**Proof.** (1): We know that $H \cap N$ is a subgroup of $H$ and it is a normal subgroup in $H$, since $N$ is normal in $G$ and $H$ is normal in $H$.

(2): As $aN = Na$ for all $a \in G$ the sets $HN$ and $NH$ are equal, so we need only show it is a group. Let $h_1, h_2 \in H$ and $n_1, n_2 \in N$. We must show $(h_1n_1)(h_2n_2)^{-1}$ lies in $HN$. But, by Great Trick,

$$(h_1n_1)(h_2n_2)^{-1} = h_1n_1n_2^{-1}h_2^{-1} = (h_1n_1n_2^{-1}h_1^{-1})(h_1h_2^{-1})$$

lies in $NH = HN$.

(3) is immediate.

(4): Define

$$\varphi : H \to HN/N \text{ by } x \mapsto xN.$$ 

As $\varphi(xy) = xyN = xNyN = \varphi(x)\varphi(y)$ for all $x, y$ in $N$, the map is a homomorphism. We show $\ker \varphi = H \cap N$. If $x \in \ker \varphi \subset H$, 

then \( N = e_{HN/N} = \varphi(x) \) so \( x \in N \). Conversely, if \( x \in H \cap N \), then \( N = xN = \varphi(x) \) and \( x \in \ker \varphi \). By definition
\[
HN/N = \{hnN \mid h \in H, n \in N\} = \{hN \mid h \in H\},
\]
so \( \varphi \) is surjective with kernel \( H \cap N \). By the First Isomorphism Theorem \( \varphi \) induces an isomorphism \( \varphi : H/H \cap N \to HN/N \). 

**Exercises 13.7.**

1. Prove the properties of preimages in 13.1.
2. Prove the alternate form of the Correspondence Principle 13.3.
3. Let \( G \) be a group and \( H \) and \( K \) normal subgroups of \( G \). We say that \( G \) is the *internal direct product* of \( H \) and \( K \) if \( G = HK \) and \( H \cap K = 1 \). Show that if \( G \) is the internal direct product of \( H \) and \( K \), then \( G \cong H \times K \). Why is \( H \times K \) not the internal direct product? For which subgroups of \( H \times K \) is \( H \times K \) the internal direct product?
4. Let \( G \) be a group with normal subgroups \( N_i, i = 1, \ldots, r \). Call \( G \) the the *internal direct product* of the \( N_i \) if \( G = \langle \bigcup_{i=1}^r N_i \rangle \) and if for each \( j, j = 1, \ldots, r, H_j \cap \bigcup_{i \neq j} N_i = 1 \). Show that \( G \cong N_1 \times \cdots \times N_r \).
5. Let \( G \) be a group with normal subgroups \( N_i, i = 1, \ldots, r \) such that \( G = N_1 \cdots N_r \). Show that \( G \) is the internal direct product of the \( N_i, i = 1, \ldots, r \) if and only if for every \( x \in G \) there exist unique \( x_i \in N_i \) such that \( x = x_1 \ldots x_n \).
6. Suppose that \( G \) is a group with normal subgroups \( N_i, i = 1, \ldots, r \) such that \( |G| = |N_1| \cdots |N_r| \) and the orders of the \( N_i \) are pairwise relatively prime. Show that \( G \) is the (internal) direct product of the \( N_i, i = 1, \ldots, r \) and if \( H \) is a subgroup of \( G \), then \( H \) is the (internal) direct product of the \( H \cap N_i, i = 1, \ldots, r \).
7. Let \( N \) be a normal subgroup of a group \( G = H \times K \). Show that either \( N \) is abelian or \( N \) intersects one of \( H \) or \( K \) non-trivially.
8. Let \( N \) and \( H \) be two normal subgroups of \( G \) such that \( G = HN \). Show that there is an isomorphism \( G/(H \cap N) \cong G/H \times G/N \).
9. Let \( G \) be a finite group with normal subgroups \( H \) and \( K \) of relatively prime order. Show that the group \( HK \) is cyclic if \( H \) and \( K \) are cyclic and abelian if \( H \) and \( K \) are abelian.
10. Let \( G \) be a group with subgroups \( H \) and \( N \). We say that \( G \) is the *(internal) semidirect product* of \( H \) and \( N \), denoted by \( N \rtimes H \), if

(i) \( N \triangleleft G \)

(ii) \( G = NH \)

(iii) \( N \cap H = \{1\} \).
Let $G = N \rtimes H$. Define $\varphi : H \rightarrow \text{Aut}(N)$ by $\varphi_x(n) = xnx^{-1}$. Show $\varphi$ is a group homomorphism and $\psi : N \rtimes H \rightarrow N \rtimes \varphi H$ induced by $h \mapsto (e_N, h)$ and $h \mapsto (n, e_H)$ is a group isomorphism, where $N \rtimes \varphi H$ is the (external) semidirect product of Exercise 11.9(14). [One often identifies these groups.]

11. Prove that the quaternion group is not a semidirect product.

12. Let $G$ be a group. Recall the commutator $G'$ of $G$ is the subgroup generated by the commutators of elements of $G$, i.e., elements of the form $[x, y] := xyx^{-1}y^{-1}$. (Cf. Exercise 11.9(19).) Show

(i) $G/G'$ is abelian. $G/G'$ is called the abelianization of $G$ and denoted by $G^{ab}$.

(ii) If $N \triangleleft G$ and $G/N$ is abelian then $G' \subseteq N$. In particular, the abelianization $G^{ab} = G/G'$ of $G$ is the maximal abelian quotient of $G$.

(iii) If $G' \subseteq H \subseteq G$ then $H \triangleleft G$.

13. Let $G$ be a finite abelian group and $p$ a (positive) prime dividing the order of $G$. Show that there exists an element of order $p$ in $G$.

14. In the previous exercise, prove the same result without assuming that $G$ is abelian.

### 14. Finite Abelian Groups

In this section, we prove special cases of results that will be significantly generalized in subsequent sections. We do so in order to expose some of the ideas that shall be used later that apply in this much simpler situation.

**Lemma 14.1.** Let $G$ be an abelian group and $H_1$, $H_2$ finite subgroups of relatively prime order. Then $H_1H_2 = H_2H_1$ is a group and if further $H_1 \cap H_2 = 1$, then $|H_1H_2| = |H_1||H_2|$. Moreover, if the groups $H_1$ and $H_2$ are both cyclic, then so is $H_1H_2$.

We leave the proof of this as an exercise. The next result is a special case of Cauchy’s Theorem (cf. Theorem 21.23 below) when $G$ is abelian.

**Proposition 14.2.** Let $G$ be a finite abelian group and $p > 0$ a prime dividing the order of $G$. Then there exists an element of $G$ of order $p$.

**Proof.** We prove this by induction on $|G|$. We may assume that $G \neq 1$. As $G$ has no nontrivial subgroups if and only if $G \cong \mathbb{Z}/p\mathbb{Z}$ for some prime $p$, we may assume that $|G|$ is not a prime. In particular, $G$ has a subgroup $1 < H < G$. If $p \mid |H|$, we are done by induction on $|G|$, so we may assume $p \nmid |H|$. Therefore, there exists a prime $q \neq p$
such that \( q \mid |H| \). By induction on \( |G| \), there exists an element \( y \) in \( H \)
of order \( q \). As \( G \) is abelian, \( \overline{G} = G/(y) \) is a group. Let \( \phi : G \to \overline{G} \) be
the canonical epimorphism. Since \( |\overline{G}| < |G| \) and \( p \mid |\overline{G}| \) by Lagrange’s
Theorem, there exists an element \( z \) in \( G \) such that \( \phi(z) \) has order \( p \) in \( \overline{G} \).
It follows that \( z^p \) lies in the kernel of \( \phi : G \to \overline{G} \), so \( z^p = y^i \) for some \( i \), hence \( z^q \) has order \( p \) in \( G \).

We now use the proposition above to prove Sylow’s First Theorem
for the case when a group is abelian. [Cf. Theorem 22.1 below.]

**Theorem 14.3.** Let \( G \) be a finite abelian group and \( p > 0 \) a prime
dividing \( |G| \), say \( |G| = p^n m \) with \( p \) and \( m \) relatively prime. Then
\[
G(p) := \{ x \in G \mid x^{p^r} = e \text{ some integer } r \} < G \text{ and } |G(p)| = p^n.
\]
Moreover, \( G(p) \) is the unique subgroup of \( G \) of order \( p^n \).

**Remark 14.4.** In the theorem, if \( G \) is finite and not abelian, there
may be more than one group of order \( p^n \). Indeed, we shall see there is
only one if and only if there is a normal subgroup of order \( p^n \).

**Proof.** The set \( G(p) \) is a subgroup of \( G \) (why?) and normal as \( G \) is
abelian. As every element in \( G(p) \) has order a power of \( p \), it follows by
the previous proposition that \( |G(p)| = p^r \) some \( r \geq 1 \) and by Lagrange’s
Theorem that \( r \leq n \). In particular, if \( x \) lies in \( G(p) \), then \( x^{p^n} = e \). We
must show that \( r = n \). Suppose to the contrary that we have \( r < n \).
Set \( \overline{G} = G/G(p) \) and let \( \phi : G \to \overline{G} \) be the canonical epimorphism.
By the Correspondence Principle (Alternative Form), \( p \mid |\overline{G}| \), hence by
the previous proposition, there exists an element \( x \in G \setminus G(p) \) and an
integer \( m > 0 \) such that
\[
\overline{x^{p^m}} = \overline{e}, \text{ i.e., } x^{p^m} \in G(p).
\]
As \( x^{p^{n+m}} = (x^{p^m})^{p^n} = e \), we have \( x \in G(p) \), a contradiction. Thus
\( |G(p)| = p^n \) as desired.

If \( H \) is subgroup of \( G \) of order \( p^n \), then \( HG(p) \) is a subgroup of
\( G \) of order \( |H||G(p)|/|H \cap G(p)| \geq p^n \) by the Second Isomorphism
Theorem. It follows that \( |H| = |H \cap G(p)| = |G(p)| = p^n \), hence
\( H = H \cap G(p) = G(p) \).

**Corollary 14.5.** Let \( G \) be a finite abelian group of order \( n = p_1^{m_1} \cdots p_r^{m_r} \)
with positive primes \( p_1 < \cdots < p_r \) and positive integers \( m_1, \ldots, m_r \).
Then \( G = G(p_1) \cdots G(p_r) \). Moreover, \( G \cong G(p_1) \times \cdots \times G(p_r) \).

**Proof.** As \( G(p_1) \cdots G(p_r) \) is a group and we know that
\[
|G(p_1) \cdots G(p_r)| = |G(p_1)| \cdots |G(p_r)|
\]
by Lemma 14.1, the first statement follows by the theorem. The map
\[ G(p_1) \times \cdots \times G(p_r) \rightarrow G \]
given by \((x_1, \ldots, x_r) \mapsto x_1 \cdots x_r\)
is a group homomorphism as \(G(p_i)G(p_j) = G(p_j)G(p_i)\) for all \(i\) and \(j\). Suppose that \(x_1 \cdots x_r = e\) with \(x_i \in G(p_i)\) for \(i = 1, \ldots, r\). As \(x_j^{n/p^m_i} = e\) for all \(j \neq i\) and \(p\) divides the order of \(x_i\), it follows that \(x_i = e\) for all \(i\). Thus the map is a monomorphism and by counting an isomorphism.

It follows, to determine all finite abelian groups up to isomorphism, we have reduced to the study of groups of order a power of a prime. If \(p > 1\) is a prime, a nontrival group of order a power of \(p\) is called a \(p\)-group.

For notational reasons, it is easier to write abelian groups additively, i.e., use additive groups. If \(H\) and \(K\) are subgroups of an additive group, we know that \(H + K\) is a subgroup. If \(H\) and \(K\) also satisfy \(H \cap K = 0\), we shall write \(H \oplus K\) for the group \(H + K\). Of course, if \(H\) and \(K\) are finite groups of relatively prime order, we know that \(H + K = H \oplus K\). The interesting case is when this is not true.

**Lemma 14.6.** Let \(G\) be a finite additive \(p\)-group and suppose that the element \(x\) in \(G\) has maximal order. Then there exists a subgroup \(H\) of \(G\) satisfying \(G = \langle x \rangle \oplus H\).

**Proof.** Let \(p^n\) be the order of \(x\). By the Well-ordering Principle, there exists a maximal subgroup \(H\) of \(G\) satisfying \(H \cap \langle x \rangle = \{0\}\). Therefore, \(\langle x \rangle + H = \langle x \rangle \oplus H\), and we are done if \(\langle x \rangle + H = G\). So suppose not. Let \(\overline{G} = G/(\langle x \rangle + H)\) and \(- : G \rightarrow \overline{G}\) the canonical epimorphism. By Proposition 14.2, there exists an element \(y \in G\) satisfying \(y \notin \langle x \rangle + H\) and \(\overline{y}\) has order \(p\) in \(\overline{G}\). In particular, \(py\) lies in \(\langle x \rangle + H\). Write
\[ py = lx + h \quad \text{with} \quad l \geq 0 \quad \text{in} \quad \mathbb{Z} \quad \text{and} \quad h \in H. \]

It follows that \(p^ny = p^{n-1}lx + p^{n-1}h\). Since \(x\) has maximal order \(p^n\) among all elements in \(G\), we have \(0 = p^ny = p^{n-1}lx + p^{n-1}h\). Therefore, \(p^{n-1}lx = -p^{n-1}h\) and it lies in \(\langle x \rangle \cap H = 0\). It follows that \(p | l\), say \(l = l'p\). Thus (*) becomes
\[ py = l'px + h \quad \text{equivalently} \quad p(y - l'x) = h \quad \text{in} \quad H. \]

As \(y \notin \langle x \rangle + H\), we must have \(y - l'x \notin \langle x \rangle + H\). The maximality of \(H\) now yields
\[ \langle y - l'x, H \rangle \cap \langle x \rangle \neq 0. \]
Therefore, there exist integers $r, s$ with $r$ nonzero satisfying
\[ 0 \neq rx = s(y - l'x) + h' \text{ for some } h' \in H. \]
Hence $sy = rx + sl'x - h'$ lies in $\langle x \rangle + H$. We look at the integer $s$.
If $p \mid s$, then by (†), we have $s(y - l'x)$ lies in $H$. This implies that $rx$
lies in $\langle x \rangle \cap H = 0$, a contradiction.
If $p \not\mid s$, then $p$ and $s$ are relatively prime, so there exists an equation
$1 = pa + sb$ for some integers $a$ and $b$. Since both $py$ and $sy$
lie in $\langle x \rangle + H$, we have $y = apy + bsy$ lies in $\langle x \rangle + H$, a contradiction. The
lemma now follows. \qed

**Corollary 14.7.** Let $G$ be a finite $p$-group. Then $G$ is a product of
cyclic $p$-groups.

**Proof.** We may assume that $G$ is additive and not cyclic. In
particular $|G| > p$. By the lemma, $G = \langle x \rangle \oplus H$ for some $x$ in $G$ and
subgroup $H$ of $G$. As $G$ is not cyclic, $H \neq G$, and the result follows by
induction on $|G|$.

**Remark 14.8.** Let $p > 0$ be a prime, $N$ a positive integer. Suppose
that $G$ is an additive group in which every element $a$ in $G$ satisfies
$p^N a = 0$. Let $x$ be an element of $G$ of maximal order. Then there exists
a subgroup $H$ of $G$ satisfying $G = \langle x \rangle \oplus H$. The proof is essentially
the same as that above except that one replaces the Well-Ordering
Principle by Zorn’s Lemma, an axiom that we shall study later. [Cf.
§27.5.]

**Proposition 14.9.** Every finite abelian group is a product of cyclic
groups.

**Proof.** By 14.5 every finite abelian group is a product of finite
abelian $p$-groups. By Corollary 14.5, every finite abelian $p$-group is a
product of cyclic $p$-groups. \qed

**Theorem 14.10.** (Fundamental Theorem of Finite Abelian Groups)
Let $G$ be a finite additive group and for each prime $p > 0$ dividing $G$,
let $G(p)$ be the unique $p$-subgroup of $G$ of maximal order. Then
\[ G = \bigoplus_{p \mid |G|} G(p) \]
Moreover, if $p \mid |G|$, then
\[ G(p) \cong \times_{i=1}^{r} \Z/p^{n_i} \Z \]
with $r$ unique and $1 \leq n_1 \leq \cdots \leq n_r$ also unique relative to this
ordering. In particular, any finite abelian group is a product of cyclic
$p$-groups for various primes $p$. 
Proof. Let \( p \mid |G| \). By Corollary 14.5 and Proposition 14.9, it suffices to show \( G(p) \cong \bigtimes_{i=1}^{r} \mathbb{Z}/p^{n_i} \mathbb{Z} \) and uniquely up to isomorphism. As every abelian \( p \)-group is isomorphic to a product of cyclic \( p \)-groups by Lemma 14.6 and every cyclic \( p \)-group must be isomorphic to \( \mathbb{Z}/p^{a_1} \mathbb{Z} \) for some integer \( a_1 \), to finish, it suffices to show that if \((*)\)

\[
\bigtimes_{i=1}^{r} \mathbb{Z}/p^{n_i} \mathbb{Z} \cong \bigtimes_{j=1}^{s} \mathbb{Z}/p^{m_j} \mathbb{Z}
\]

with \( n_1 \geq \cdots \geq n_r \) and \( m_1 \geq \cdots \geq m_s \) (we changed notation for convenience), then \( r = s \) and \( n_i = m_i \) for all \( i \). We may assume that the number \( N \) of the \( n_i \) that are equal to 1 is less than the number \( M \) of the \( m_j \) that equal 1. In particular, \( p^N \leq p^M \). Multiplying \((*)\) by \( p \), we see that

\[
\bigtimes_{i=N+1}^{r} \mathbb{Z}/p^{n_i-1} \mathbb{Z} \cong \bigtimes_{j=M+1}^{s} \mathbb{Z}/p^{m_j-1} \mathbb{Z}
\]

as \( p(\mathbb{Z}/p^k \mathbb{Z}) \cong \mathbb{Z}/p^{k-1} \mathbb{Z} \) for all \( k \). By induction, \( n_i - 1 = m_i - 1 \) for all \( i > N \). As \( |G(p)| = \prod_{i=1}^{r} p^{n_i} = \prod_{i=1}^{s} p^{m_i} \), we conclude that \( r = s \) and \( n_i = m_i \) for all \( i \). \( \square \)

Remarks 14.11. The theorem generalizes to finitely generated abelian groups. We indicate how this would proceed without a full proof. An abelian group \( G \) is called

(i) torsion-free if every non-identity element in \( G \) has infinite order.
(ii) torsion if every non-identity element in \( G \) has finite order.

Let \( G \) be an abelian group and set

\[ G_t := \{ x \in G \mid x \text{ has finite order} \}. \]

By the Binomial Theorem, \( G_t \subset G \) is a subgroup and by the Correspondence Principle, \( G/G_t \) is torsion-free. If \( G \) is finitely generated, then one can show that \( G_t \) is a finitely generate torsion abelian group hence finite, so decomposes uniquely as in the theorem and, moreover, that \( G \cong G_t \times G/G_t \). One also shows if \( G \) is a finitely generated torsion-free abelian group, then \( G \cong \mathbb{Z}^r \) for a unique integer \( r \). Therefore, any finitely generated abelian group \( G \) is isomorphic to \( \bigtimes_{p \mid |G|} G(p) \times \mathbb{Z}^r \) for some unique \( r \) and the \( G(p) \) have unique decompositions given by Theorem 14.10. We shall prove this and an important generalization of this that has significant application to the theory of matrices in \( \S 43 \).

Exercises 14.12.

1. Prove Lemma 14.1. If we assume that \( G \) is an arbitrary group, what conditions on \( H_1 \) and \( H_2 \) will still guarantee the result. Prove this.
2. Let \( G \) be an abelian group. Show that \( G(p) := \{ x \in G \mid x^{p^r} = e \text{ some } r \} \) is a subgroup of \( G \).
3. Show that \( p(\mathbb{Z}/p^k\mathbb{Z}) \cong \mathbb{Z}/p^{k-1}\mathbb{Z} \) for all primes \( p \) and positive integers \( k \).

4. Determine up to isomorphism all finite abelian groups up to order 30.

5. Determine all \( n > 1 \) such that all abelian groups of order \( n \) are cyclic and prove your determination.

15. Addendum: Divisible Groups

In this section, we investigate another class of abelian groups consisting of non-finitely generated abelian groups (if we exclude the trivial group). These abelian groups are called divisible groups. They arise naturally from the divisibility property of the rational numbers. It is convenient to often use additive notation for our abelian group, and will call them additive groups as usual when we do so. In this section, we shall also need to use some results to be proven in later sections, the major one is a proof of Theorem 15.13 (which requires Zorn’s Lemma) as does the fact that all vector spaces have bases. We shall also leave many details as exercises.

When discussing sets of infinitely many abelian groups, we must distinguish between two group constructions. Let \( G_i, i \in I \), be groups.

We have defined the direct product \( \prod_{i \in I} G_i := \{(g_i)_{i \in I} \mid g_i \in G_i, i \in I\} \). If all the groups \( G_i, i \in I \), are additive groups, the direct product contains a proper subgroup called the (external direct sum) of the \( G_i \)'s defined by

\[
\prod_{i \in I} G_i := \{(g_i)_{i \in I} \in \prod_{i \in I} G_i \mid g_i = e_{G_i} \text{ for all but finitely many } i \in I\}.
\]

As in Section 14, we also want to look at an abelian group \( G \) containing subgroups \( G_i, i \in I \), with

\[
G = \sum_{i \in I} G_i := \{\sum_{i \in I} g_i \mid g_i \in G_i \text{ and } g_i \neq e_G \text{ for finitely many } i \in I\}
\]

that is isomorphic to \( \prod_{i \in I} G_i \). The condition for this to be true is that if

\[
\sum_{i \in I} g_i = 0 \text{ in } G, g_i \in G_i, i \in I, \text{ then } g_i = e_G \text{ for all } i \in I.
\]

(Cf. this to linear independence in vector spaces.) If this is true, then we write \( G = \bigoplus_{i \in I} G_i \) and call it the (internal) direct sum of \( G \) by the \( G_i \)'s. (Cf. Exercise 15.18(1).)
We extend our definition of the subgroup $G(p)$, $p$ a prime, of a finite additive group $G$ used in Section 14 to: If $G$ is an additive group and $p > 0$ a prime, let
\[ G(p) := \{ x \in M \mid p^r x = 0 \text{ for some positive integer } r \}, \]
a torsion group.

**Definition 15.1.** Let $\mathcal{P}$ denote the positive primes. Let $G$ be an additive group and $p \in \mathcal{P}$. We call $G(p)$ the $p$-primary part of $G$. If $G = G(p)$, we call $G$ a $p$-primary group. We also let
\[ pG := \{ g \in G \mid pg = 0 \}, \]
the subgroup generated by elements of order $p$ in $G$ (or $G(p)$).

We shall also need a suitable generalization of Fundamental Theorem Finite Abelian Groups 14.10. The general version is called the Primary Decomposition Theorem that we shall later prove in Theorem 43.20 below. The case that we need below in Theorem 15.2 can be proven by facts that you know (and Exercise 15.18(1)), so we omit it.

**Theorem 15.2.** Let $G$ be a torsion additive group. Then $G = \bigoplus_{p \in \mathcal{P}} G(p)$.

If $G$ is an abelian group, it was mentioned in Section 14 [to be later in Corollary 43.16], that $G \cong G_t \times (G/G_t)$ if $G$ is finitely generated. It is not true in general, however, that $G \cong G_t \bigoplus (G/G_t)$ when $G$ is not finitely generated. We give a counterexample.

**Example 15.3.** Let $G = \prod_{p \in \mathcal{P}} \mathbb{Z}/p\mathbb{Z}$ as an additive group. We show that $G \not\cong G_t \prod (G/G_t)$.

Let $x = (x_p)_{p \in \mathcal{P}} \in G$. Suppose that $q \in \mathcal{P}$ satisfies $qy = x$ for some $y \in G$, i.e., $qy_p = x_p$ for all $p \in \mathcal{P}$. We say that $x$ is divisible by $q$. Then $x_q = 0$. In particular, the only element $x \in G$ divisible by every $p \in \mathcal{P}$ is $x = 0$, i.e., that there cannot be a nonzero element in $G$ that is divisible by every prime $p \in \mathcal{P}$.

Let $\gamma : G \to G/G_t$ be the natural group epimorphism. We show that $G/G_t$ contains a nonzero element $\overline{a} \in G/G_t$, $a \in G$, divisible by all primes $p \in \mathcal{P}$. Let $a_p$ be a generator of $\mathbb{Z}/p\mathbb{Z}$ and $a = (a_p)_{p \in \mathcal{P}}$. Fix $q \in \mathcal{P}$. Then for every prime $p \neq q$, there exists $x_p \in \mathbb{Z}/p\mathbb{Z}$ satisfying $qx_p = a_p$. Set $x_q = 0$ and let $y \in G$ be the element with entries all 0 except for the $q$th coordinate that is $a_q$. In particular, $y \in G_t$ and $qx = a - y$. It follows that $q\overline{x} = \overline{a} - \overline{y} = \overline{a}$. Since this is true for every prime $q$, the element $\overline{a}$ in $G/G_t$ is divisible by every prime. Since $a$ cannot have finite order in $G$, we have $\overline{a}$ is nonzero and the assertion is verified. In particular, if $G \cong G_t \prod (G/G_t)$, then there exists a nonzero
element in $G$ divisible by all primes, contradicting that no such element lies in $G$.

We shall show that that we do have such a decomposition for one type of abelian group. To do so, we will need a few facts about vector spaces. We shall assume that all vector spaces have bases. This will be proven in Proposition 27.6 below. By linear algebra, we then have two vector spaces over a field $F$ are isomorphic if and only if they have bases of the same cardinality. For any $F$-vector space $V$, we write $\dim V$ for the cardinality of a basis (it is well-defined). We leave the proofs of the following three lemmas as exercises.

**Lemma 15.4.** Let $V$ be a vector space over a field $F$. Suppose that $V$ has a basis $B = \{v_i\}_I$ Then $(V,+)=\bigoplus_B Fx$ as an additive group.

**Lemma 15.5.** Let $V$ and $W$ be vector spaces over $F$ with $F$ either $\mathbb{Q}$ or $\mathbb{Z}/p\mathbb{Z}$ with $p > 0$ a prime. Then $V$ and $W$ are isomorphic as vector spaces if and only if they are isomorphic as additive groups.

Lemma 15.5 is false in general, e.g., it can be shown $\mathbb{R} \ncong \mathbb{R}^2$ as abelian groups using that they have bases as vector spaces over $\mathbb{Q}$ of the same uncountable cardinality.

**Lemma 15.6.** Let $G$ be an additive group and $p > 0$ a prime. Then $pG$ is a vector space over $\mathbb{Z}/p\mathbb{Z}$.

**Definition 15.7.** Let $G$ be an additive group. We say that $G$ is a **divisible group** if for all $y \in G$, $n \in \mathbb{Z}^+$, there exists an $x \in G$ satisfying $nx = y$, i.e., $\frac{1}{n}y \in G$. Equivalently, the group homomorphism $G \to G$ given by $x \mapsto nx$ is surjective for all positive integers $n$.

**Examples 15.8.**
1. $\mathbb{Q}$, $\mathbb{R}$, $\mathbb{C}$ are all torsion-free divisible additive groups.
2. $\mathbb{R}^+$ is a multiplicative torsion-free divisible group
3. If $G$ is a torsion-free divisible group, then $G \to G$ by $x \mapsto nx$ is an isomorphism for all $n \in \mathbb{Z}^+$.
4. $\mathbb{C}^\times$ is a divisible group. The circle group $T$ in $\mathbb{C}^\times$ is the torsion subgroup of $\mathbb{C}^\times$.

Checking that the axioms of a vector space hold, we see that the following is true.

**Lemma 15.9.** A torsion-free additive divisible group is a vector space over $\mathbb{Q}$.

**Corollary 15.10.** Let $G$ be a torsion-free divisible group. Then $G \cong \prod_I \mathbb{Q}$ for some indexing set $I$. 
The following observations are easily seen:

**Observations 15.11.** Let $G$ be a nonzero divisible group. Then
1. $G$ is an infinite group.
2. $G_t$ is a divisible group.
3. If $\varphi : G \to G'$ is a group epimorphism, then $G'$ is also a divisible group. In particular, $G/G_t$ is a torsion-free divisible group.
4. Let $G_i, i \in I$, be additive groups. Then $G_i$ is a divisible group for all $i \in I$ if and only if $\prod_I G_i$ is a divisible group if and only if $\bigoplus_I G_i$ is a divisible group.

**Example 15.12.** By Observation 15.11(3), the additive group $\mathbb{Q}/\mathbb{Z}$ is a divisible torsion group. It is not a finitely generated group, hence neither is $\mathbb{Q}$. Let $p$ be a prime number. Define $\mathbb{Z}_{p^\infty} := (\mathbb{Q}/\mathbb{Z})(p)$. Then $\mathbb{Z}_{p^\infty}$ is a divisible group and is not finitely generated. (Why?) Moreover, by Theorem 15.2, we have $\mathbb{Q}/\mathbb{Z} = \bigoplus_p \mathbb{Z}_{p^\infty}$.

We wish to classify divisible groups. To do so, we shall need a defining property of divisible groups. Unfortunately, to prove this theorem, we need to use Zorn’s Lemma to be done later, so we shall postpone its proof until then [cf. Proposition 27.9]. We shall, however, assume its validity here.

**Theorem 15.13.** Let $G$ be a divisible abelian group and $B$ an abelian group with $A \subset B$ a subgroup. Suppose that $\varphi : A \to G$ is a group homomorphism. Then there exists a group homomorphism $\psi : B \to D$ satisfying $\psi|_A = \varphi$. We say that $\psi$ extends $\varphi$ to $G$.

**Corollary 15.14.** Let $G$ be an abelian group and $H$ a subgroup of $G$ that is a divisible group. Then $G = H \oplus X$, for some subgroup of $G$.

**Proof.** The identity map $1_H : H \to H$ extends to $\psi : G \to H$, i.e., $\psi|_H = 1_H$ by the theorem. We show that $G = H \oplus \ker \psi$. If $x \in G$, then $\psi(x) \in H$, hence $\psi(x - \psi(x)) = \psi(x) - 1_H \psi(x) = 0$. Therefore, $(x - \psi(x)) + \psi(X)$ lies in $H + \ker \psi$. If $x \in H \cap \ker \psi$, then $x = \psi(x) = 0$. This shows that $G = \ker \psi \oplus H$ by Exercise 15.18(1) as needed.

Using Theorem 15.13, we prove the following lemma.

**Lemma 15.15.** Suppose that $G_1$ and $G_2$ are divisible $p$-primary groups. Then $G_1 \cong G_2$ if and only if $pG_1 = pG_2$.

**Proof.** ($\Rightarrow$) follows easily from Lemma 15.6
(⇐): Let \( \varphi : pG_1 \to pG_2 \) be a group isomorphism. We view it as a monomorphism into \( G_2 \). By Theorem 15.13, we can extend \( \varphi \) to \( \psi : G_1 \to G_2 \), i.e., \( \psi|_{pG_1} = \varphi \). We show that \( \psi \) is an isomorphism.

\( \psi \) is injective: We show that \( x = 0 \) by induction on elements of order \( p^r \) in \( G_1 \). Let \( x \) be a nonzero element in \( G_1 \) of order \( p^r \) and suppose that \( \psi(x) = 0 \). If \( r = 1 \), then \( px = 0 \) and \( x = 0 \) as \( \varphi \) is an isomorphism. Let \( r > 1 \). Since \( px \) has order \( p^{r-1} \), by induction \( px = 0 \), contradicting \( x \) has order \( p^r \). Therefore, \( \psi \) is injective.

\( \psi \) is surjective: Let \( y \in G_2 \) have order \( p^r \). If \( r = 1 \), since \( \varphi \) is an isomorphism, there exists an \( x \in G_1 \) such that \( \psi(x) = y \). So assume that \( r > 1 \). Since \( 0 \neq p^{r-1}y \in pG_2 \), there exists \( x \in pG_1 \) such that \( \psi(x) = p^{r-1}y \). Since \( G_1 \) is divisible, there exist \( z \in G_1 \) such that \( p^{r-1}z = x \). Therefore, \( p^{r-1}(y - \psi(z)) = 0 \). By induction on \( r \), there exists an element \( w \in G_1 \) satisfying \( \psi(w) = y - \psi(z) \). Therefore, \( y = \psi(w + z) \). \( \square \)

We can now classify divisible abelian groups.

**Theorem 15.16.** Let \( G \) be a divisible additive group. Then

\[
G = \bigoplus_{I} \mathbb{Q} \bigoplus \left( \bigoplus_{p} \bigoplus_{I_p} \mathbb{Z}_{p^\infty} \right)
\]

for some indexing sets \( I, I_p, p \in \mathcal{P} \) (possibly empty).

**Proof.** Since \( G_t \) is divisible by Observation 15.11(2), we have \( G \cong G_t \bigoplus (G/G_t) \) by Corollary 15.14. Since \( G/G_t \) is torsion-free, it is a direct sum of copies of \( \mathbb{Q} \) by Lemma 15.10. By Observation 15.11(4) and Theorem 15.2 each \( G(p) \) is divisible. Finally, let \( G'_p \) be the direct sum of \( \dim_{\mathbb{Z}/p\mathbb{Z}}(G) \) copies by of \( \mathbb{Z}_{p^\infty} \). By Lemma 15.15, the groups \( G(p) \) and \( G'_p \) are isomorphic. \( \square \)

Divisible groups have the important property that given any abelian group \( A \), there exits a divisible abelian group \( G \) and a group monomorphism \( \varphi \) : \( A \to G \). To establish this, we need some further ideas from vector space theory.

The additive group, the direct sum \( \bigoplus_I \mathbb{Z} \) has a \( \mathbb{Z} \)-basis (with scalars from \( \mathbb{Z} \)) just as with an \( F \)-vector space over a field \( F \) (with scalars from \( F \)). In fact, \( S := \{ e_i \mid i \in I \} \) is a basis for \( \bigoplus_I \mathbb{Z} \) where \( e_i \) has 1 in the \( i \)th coordinate and 0 in the \( j \)th coordinate for all \( j \neq i \) in \( I \). So every element \( x \) is a finite sum of some of the \( e_i \), i.e., \( x = \sum_I a_i e_i \) with \( a_i \in \mathbb{Z} \), all but finitely many \( a_i \) nonzero. This means that \( S \) spans. In addition, the \( a_i \) are unique as \( y = 0 \) in \( \bigoplus_I \mathbb{Z} \) if and only all the coordinates are zero. This means that \( S \) is linear independence.
Just as for vector spaces, this means that if $G$ is an abelian group and $a_i \in \mathbb{Z}$, $i \in I$, (not necessarily distinct), there exists a unique group homomorphism induced by $e_i \mapsto a_i$ for all $i \in I$. [The same proof used for vector spaces works.] For example, since $S \subset \coprod I \mathbb{Q}$ (and is a $\mathbb{Q}$-basis for it), the inclusion map $\text{inc} : \coprod I \mathbb{Z} \to \coprod I \mathbb{Q}$ via $e_i \to e_i$ is a group monomorphism.

**Theorem 15.17.** Let $A$ be an abelian group. Then there exist a divisible abelian group $G$ and a group monomorphism $\varphi : A \to G$.

**Proof.** Let $\{a_i \mid i \in I\}$ be a set of generators for $A$. Then there exists a group homomorphism $\psi : \coprod I \mathbb{Z} \to A$ determined by $e_i \mapsto a_i$ for all $i \in I$. Let $K = \ker \psi$. By the First Isomorphism Theorem, $\psi$ induces a group isomorphism $\overline{\psi} : (\coprod I \mathbb{Z})/K \to A$. The group monomorphism $\text{inc} : \coprod I \mathbb{Z} \to \coprod I \mathbb{Q}$ takes $K \to \mathbb{Q}K$, where $\mathbb{Q}K$ is the $\mathbb{Q}$-vector space spanned by $\{e_i \mid i \in I\}$. Since $\mathbb{Q}K \cap \coprod I \mathbb{Z} = K$ (why?), the map $(\coprod I \mathbb{Z})/K \to (\coprod I \mathbb{Q})/K$ is a group monomorphism. By Observation 4(3), the group $(\coprod I \mathbb{Q})/K$ is a divisible group. Therefore, the composition $\varphi := \text{inc} \circ \overline{\psi} : A \to (\coprod I \mathbb{Q})/K$ works. \qed

**Exercises 15.18.**
1. Let $G$ be an abelian group and $G_i$, $i \in I$ subgroups satisfying $G = \sum_i G_i$. Prove that the following are equivalent:
   (i) $G = \bigoplus_i G_i$.
   (ii) If $\sum_i g_i = 0$ in $G$, then $g_i = e_G$ for all $i \in I$.
   (iii) For all $j \in I$, we have $G_j \cap \sum_{I \setminus \{j\}} G_i = \{e_G\}$.
   (iv) $\prod_i G_i \to \sum_i G_i$ given by $(g_i)_I \mapsto \sum_i g_i$ is a group isomorphism.
2. Prove Lemma 15.4
3. Prove Lemma 15.5
4. Prove Lemma 15.6
5. Prove Theorem 15.2
6. Prove Lemma 15.9
7. Let $G$ be the group in Example 15.3. Show that $G/G_i$ is a divisible group.
8. Let $p$ be a prime number. Show all of the following:
   (i) $\mathbb{Z}_{p^\infty}$ is a divisible group.
   (ii) The set $\{\frac{1}{p^r} + \mathbb{Z} \mid r \text{ a non-negative integer } \}$ generates $\mathbb{Z}_{p^\infty}$. In particular, $\mathbb{Z}_{p^\infty}$ is not finitely generated.
(iii) The subgroup $\left< \frac{1}{p^r} + \mathbb{Z} \right>$ of $\mathbb{Z}_{p^\infty}$ is isomorphic to $\mathbb{Z}/p^r\mathbb{Z}$ for all $r \in \mathbb{Z}^+$.

9. Prove Observations 15.11.

10. Let $G_1$ and $G_2$ be divisible abelian groups. Prove that $G_1 \cong G_2$ if and only if $\dim_{\mathbb{Q}} (G_1/(G_1)_t) = \dim_{\mathbb{Q}} (G_2/(G_2)_t)$ and $\dim_{\mathbb{Z}/p\mathbb{Z}} (pG_1) = \dim_{\mathbb{Z}/p\mathbb{Z}} (pG_2)$ for all primes $p$.

11. Prove that the following groups are isomorphic: $\mathbb{R}/\mathbb{Z}$, the circle group $T$, $\prod_p \mathbb{Z}_{p^\infty}$, $\mathbb{Q}\mathbb{Z} \oplus \mathbb{R}$.

12. An abelian group $G$ is divisible if and only if it satisfies the property in Theorem 15.13.

13. An abelian group $G$ is divisible if and only if every nonzero quotient of $G$ is infinite.

14. An abelian group $G$ is divisible if and only if it has no maximal subgroups, i.e., a subgroup $M < G$ so that there exists no subgroup $H$ in $G$ with $M < H < G$.

16. Addendum: Finitely Generated Groups

Most of what we do in the study of group theory in these pages involves finite groups. In this short section, we show how the General Cayley Theorem 12.4 can be used to obtain a few results about finitely generated groups, i.e., a group $G$ containing a finite set $S$ of elements satisfying $G = \langle S \rangle$. Of course, finite groups are finitely generated, but the results that we obtain here are trivial for finite groups, as the questions we ask are about finitely generated groups that contain a proper subgroup of finite index.

We have seen that if $G$ is an arbitrary group and $H$ a subgroup of finite index $n$ in $G$, then the General Cayley Theorem gives a group homomorphism

$$\lambda : G \to \Sigma(G/H) \cong S_n$$
given by $x \mapsto \lambda_x : gH \mapsto xgH$,

with $\ker \lambda \subset H \subset G$ (and the largest normal subgroup of $G$ in $H$). Hence, by the First Isomorphism Theorem, $\lambda$ induces a group monomorphism

$$\overline{\lambda} : G/\ker \lambda \to \Sigma(G/H) \cong S_n$$
given by $x \mapsto \lambda_x : g \mapsto gH$.

Since $\Sigma(G/H)$ is a finite group, so is $G/\ker \lambda$, i.e., $\ker \lambda \subset G$ has finite index. Therefore, $G$ contains a normal subgroup of finite index. [Cf. Exercise 12.11(6).] We want to strengthen this result if $G$ is finitely generated.
Proposition 16.1. Let $G$ be a finitely generated group and $n$ a positive integer. Then there exist finitely many subgroups (if any) of $G$ of index $n$.

PROOF. Let $G = \langle a_1, \ldots, a_r \rangle$ and

(*) $\varphi : G \rightarrow S_n$

be a group homomorphism. Then the map $\varphi$ is completely determined by

$\varphi(a_1), \ldots, \varphi(a_r)$

in the finite set $S_n$. Thus there exist only finitely many possible group homomorphisms $\varphi$ in (*), hence finitely many normal subgroups $N$ of $G$ of finite index with $N = \ker \varphi$ for such $\varphi$. Fix a $\varphi$. By the First Isomorphism Theorem, we have $G/\ker \varphi \cong \varphi(G) \subset S_n$ is a subgroup. As $S_n$ is finite, the Correspondence Principle implies that there exist finitely many subgroups $H$ of $G$ satisfying $\ker \varphi \subset H \subset G$, in particular, there can be only finitely many $H$ (if any) with $[G : H] = n$. As there are finitely many $\varphi$, the result follows.

Our strengthening of the goal above can now be stated and established. Recall that a subgroup $K$ of a group $G$ is called characteristic if $\sigma|_K$ is an automorphism of $K$ (equivalently, $\sigma(K) = K$) for every automorphism $\sigma$ of $G$.

Corollary 16.2. Let $G$ be a finitely generated group. Suppose $G$ contains a subgroup $H$ of finite index. Then there exists a characteristic subgroup $K$ of $H$ of finite index in $G$.

PROOF. By the proposition, there exist finitely many subgroups

$H = H_1, \ldots, H_m$

of finite index $[G : H]$. It follows, given an automorphism $\varphi$ of $G$, that $\varphi(H) = H_i$ for some $i$, i.e., $\varphi$ permutes the $H_i$. Let $K := \cap_{i=1}^m H_i$. Then we have $\varphi(K) = K$, so $K$ is a characteristic subgroup of $H$. By Poincaré’s Lemma (cf. Exercise 10.16(6)), $K$ is of finite index in $G$.

The second result is much harder to prove.

Theorem 16.3. Let $G$ be a finitely generated group and $H$ be a subgroup of finite index. Then $H$ is a finitely generated group.

PROOF. Let $G = \langle a_1, \ldots, a_r \rangle$ and

$y_1 = e, \ldots, y_n$
be a transversal (i.e., a system of representatives for the cosets) of $H$ in $G$. Let

$$\lambda : G \to \Sigma(G/H) \text{ be given by } a \mapsto \lambda_x : aH \mapsto xaH.$$ 

Therefore, $\lambda_x$ permutes $y_1H, \ldots, y_nH$. Fix $i, 1 \leq i \leq r$. Then for each $j, 1 \leq j \leq n$, there exists a $k = k(i, j)$ with $1 \leq k \leq n$ satisfying

$$a_iy_jH = y_kH.$$ 

Therefore, there exist $h_{ij}$ in $H, 1 \leq i \leq r$ and $1 \leq j \leq n$ satisfying

$$y_k = a_iy_jh_{ij}.$$ 

Let $H_0 := \langle h_{ij} \mid 1 \leq i \leq r, 1 \leq j \leq n \rangle \subset G$. Then $H_0$ is clearly a finitely generated subgroup of $G$ contained in $H$. So it suffices to show that $H \subset H_0$. Set

$$W := \bigcup y_jH_0.$$ 

For each $i, 1 \leq i \leq r$, we have

$$a_iW = \bigcup_{j=1}^n a_iy_jH_0 = \bigcup_{j=1}^n y_{k(i,j)}h_{ij}H_0 = \bigcup_{k=1}^n y_kH_0 = W,$$

since

$$\{a_iy_1H, \ldots, a_iy_n\} = \{y_1, \ldots, y_n\}.$$ 

Consequently, $a_iW = W$ for each $i = 1, \ldots, r$. As $G = \langle a_1, \ldots, a_r \rangle$, we have

$$GW = W = \bigcup_{j=1}^n H_0.$$ 

Since $e \in W$, we conclude that

$$G = W = \bigcup_{j=1}^n y_jH_0 = \bigvee_{j=1}^n y_jH_0.$$ 

In particular,

$$H \subset G = \bigvee_{j=1}^n y_jH_0.$$ 

But $H = y_1H$ is disjoint from $\bigcup_{j=2}^n y_jH_0$ as it is disjoint from $\bigcup_{j=2}^n y_jH$, so we conclude that $H = H_0$ as needed.

$\square$
17. Series

One way of studying algebraic objects is to break them down into simpler pieces. For example, if $G$ is a group, it may be the direct product of groups, e.g., $\mathbb{Z}/6\mathbb{Z} \cong \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/3\mathbb{Z}$. Unfortunately, this does not happen very often for groups. An alternative approach would be to take a given non-trivial group $G$ and find the largest normal subgroup $N_1$ in $G$ satisfying $N_1 < G$. By the Correspondence Principle, $G/N_1$ then would be a simple group, and we have, in some sense reduced the study of $G$ to the groups $N_1$ and $G/N_1$. We then proceed with the same analysis on $N_1$, etc. Another approach would be to find a normal subgroup $N_1$ of $G$ such that $G/N_1$ is abelian (or cyclic), then continue on $N_1$, etc. In any of these cases, we produce a chain of subgroups $G > N_1 > N_2 > \cdots$ (where $A > B$ means $A \supset B$, $A \neq B$). Then various possibilities open:

1. This process never stops.
2. There exists an $i$ such that $N_i$ is nice, e.g., in our cases simple, abelian, cyclic, respectively.
3. At some point, there does not exist such an $N_i$.

Of course, we are interested in the second possibility.

**Definition 17.1.** Let $G$ be a non-trivial group. A sequence of groups

(*) \[ N_0 \subset N_1 \subset N_2 \subset \cdots \subset N_n = G \]

is called a **subnormal** (respectively, **normal**) series if $N_i \triangleleft N_{i+1}$ (respectively, $N_i \triangleleft G$) for all $i$, and **proper** if $N_i < N_{i+1}$ for all $i$. The quotients $(G/N_{n-1}) = N_n/N_{n-1}, \ldots, N_1/N_0$ are called the **factors** of the series. If (*) is a subnormal series for which $N_0 = 1$ and $N_n = G$, then it is called

(i) a **cyclic series** if $N_{i+1}/N_i$ is cyclic for all $i$.

(ii) an **abelian series** if $N_{i+1}/N_i$ is abelian for all $i$.

(iii) a **composition series** if $N_{i+1}/N_i$ is simple for all $i$.

A group is called **solvable** if it has an abelian series and **polycyclic** if it has a cyclic series. (Since the trivial group is cyclic, we also say that it has a cyclic series.)

**Example 17.2.** Every abelian group is solvable and every cyclic group is polycyclic.

One of the major results in finite group theory is the Feit-Thompson Theorem that every finite group of odd order is solvable. It is the first major theorem in the classification of finite simple groups. Its proof takes over 250 journal pages.
17. SERIES

Theorem 17.3. The following are true:

1. A subgroup of a solvable group is solvable.
2. The homomorphic image (i.e., the image of a group under a
   group homomorphism) of a solvable group is solvable.
3. If $N ◁G$ and both $N$ and $G/N$ are solvable then so is $G$.

Proof. We leave this as an exercise. It is important that you do
this exercise, as it will teach you how to use the isomorphism theorems. □

Let $G$ be a group. The commutator of $x$ and $y$ in $G$ is $[x, y] := xyx^{-1}y^{-1}$. The group generated by all commutators is called the derived subgroup of $G$ and denoted either by $[G, G]$ or by $G’$. By recursion, we define the $n$th derived subgroup of $G$ by $G^{(n)} = [G^{(n-1)}, G^{(n-1)}]$.

We leave the following easy computations as an exercise.

Properties 17.4. Let $G$ be a group. If $a, b, c$ are elements of $G$ and $σ : G → G_1$ a group homomorphism. Then

1. $[a, b]^{-1} = [b, a]$.
2. $σ([a, b]) = [σ(a), σ(b)]$.
3. $[a, bc] = [a, b][a, c]b^{-1} = [a, b][bab^{-1}, bcb^{-1}]$.

Proposition 17.5. Let $G$ be a group. Then $G^{(n)}$ is a characteristic subgroup of $G^{(n-1)}$ for all $n$. In particular, $G^{(n)}$ is characteristic in $G$, and

$G^{(n)} ⊂ G^{(n-1)} ⊂ ⋯ ⊂ G^{(1)} ⊂ G^{(0)} = G$

is a normal series (even a characteristic series — obvious definition).

Proof. Being characteristic is transitive, so this follows immediately, using Properties 17.4 □

Theorem 17.6. Let $G$ be a group. Then $G$ is solvable if and only if there exist an integer $n$ such that $G^{(n)} = 1$.

Proof. By Exercise 11.9(19), each factor group $G^{(i)}/G^{(i+1)}$ is abelian. Hence if $G^{(n)} = 1$ for some $n$, then $G$ is solvable. Conversely, suppose that

$1 = N_n ⊂ N_{n-1} ⊂ N_{n-2} ⊂ ⋯ ⊂ N_1 ⊂ G$

is an abelian series. It suffices to show that $G^{(i)} ⊂ N_i$ for all $i$. By Exercise 13.7(12), we know that $G’ ⊂ N_1$. By induction, we may assume that $G^{(i-1)} ⊂ N_{i-1}$. But then we have $G^{(i)} = [G^{(i-1)}, G^{(i-1)}] ⊂ [N_{i-1}, N_{i-1}] ⊂ N_i$, again using Exercise 13.7(12). □

We turn to composition series. We need the following easy lemma.
Lemma 17.7. (Dedekind’s Modular Law) Let $A$, $B$, and $C$ be three subgroups of $G$. If $A$ is a subset of $C$, then $A(B \cap C) = (AB) \cap C$.

Proof. $\subseteq$: If $a \in A$, $x \in B \cap C$, then $ax \in AB \cap AC \subseteq (AB) \cap C$.
$\supseteq$: If $a \in A$, $b \in B$ satisfy $ab \in C$, then $b = a^{-1}C \subseteq C$. □

We next prove an elaborate form of the Second Isomorphism Theorem. For notational simplicity, we use right cosets.

Lemma 17.8. (Zassenhaus Butterfly Lemma) Let $G$ be a group containing four subgroups $h$, $H$, $k$, and $K$ satisfying $h \lhd H$ and $k \lhd K$. Then we have the following:

1. $(H \cap k)(h \cap K) \lhd H \cap K$.
2. $k(h \cap K) \lhd k(H \cap K)$.
3. $h(H \cap k) \lhd h(H \cap K)$.
4. There exist isomorphisms:
   $$k(H \cap K)/k(h \cap K) \cong H \cap K/(H \cap k)(h \cap K)$$
   $$\cong h(H \cap K)/h(H \cap k).$$

We have the following picture illustrating the subgroups, where a group at the bottom of a line means that it is a subgroup of the group at the other end. We shall see the quotients (top/bottom) of like dotted, double, dashed lines, respectively, are isomorphic.
Proof. We use the Second and Third Isomorphism Theorems. As \( k \triangleleft K \) in the lemma and \( H \cap K \subset K \) is a subgroup, we have \( k(H \cap K) \) is a group satisfying \( k \triangleleft k(H \cap K), \) and \( H \cap k = (H \cap K) \cap k \triangleleft H \cap K \) together with an isomorphism

\[
(* \quad k(H \cap K)/k \overset{\sim}{\longrightarrow} H \cap K/H \cap k \quad \text{given by} \quad k a \mapsto (H \cap k)a.
\]

Similarly, as \( h \triangleleft H \) and \( H \cap K \subset H \) is a subgroup, we have \( h(H \cap K) \) is a group satisfying \( h \triangleleft h(H \cap K), \) and \( h \cap K \triangleleft H \cap K \) together with an isomorphism

\[
\quad h(H \cap K)/h \overset{\sim}{\longrightarrow} H \cap K/h \cap K \quad \text{given by} \quad ha \mapsto (h \cap K)a.
\]

As the product of normal subgroups is normal, we also have

\[
(\dagger \quad (H \cap k)(h \cap K) \triangleleft H \cap K.
\]

By restriction, the isomorphism in (*) induces an isomorphism

\[
(* \quad k(H \cap k)(h \cap K)/k \overset{\sim}{\longrightarrow} (H \cap k)(h \cap K)/H \cap k.
\]

By the Modular Law, we have \( k(H \cap k) = kH \cap k = k. \) Consequently, \( k(H \cap k)(h \cap K) = k(h \cap K). \) Applying the Correspondence Principle to the isomorphism in (*), we know that \( (H \cap k)(h \cap K) \triangleleft H \cap K \) corresponds to \( k(h \cap K) = k(H \cap k)(h \cap K) \triangleleft k(H \cap K). \)

Using the Third Isomorphism Theorem, we conclude by (*) and (*) that

\[
k(H \cap K)/k(h \cap K) \cong \frac{k(H \cap K)/k}{k(h \cap K)/k}
\]

\[
\cong \frac{H \cap K/ H \cap k}{(H \cap k)(h \cap K)/H \cap k}
\]

\[
\cong H \cap K/(H \cap k)(h \cap K).
\]

The other isomorphism is obtained similarly.

Next we wish to show if a group has a composition series, then the number of terms in a minimal, proper composition series is constant with factors unique up to isomorphism.

Definition 17.9. Given a proper subnormal series of a group \( G, \) we say another subnormal series is a refinement of the given series, if we add in new subgroups, and it is called a proper refinement if the resulting subnormal series is proper. Two proper subnormal series for \( G \) are called equivalent if they have the same number of terms, say \( n, \) and the \( n \) factors of each are isomorphic up to a permutation.
If a group has a composition series, we want to prove that the number of factors in a composition series is unique as well as these factors are unique up to isomorphism and order. The next theorem is the key to this. We use the Butterfly Lemma to prove it.

**Theorem 17.10.** (Schreier Refinement Theorem) Any two proper subnormal series for a group $G$ have equivalent refinements.

**Proof.** Suppose

\[ 1 = N_0 \triangleleft N_1 \triangleleft \cdots \triangleleft N_r = G \]
\[ 1 = H_0 \triangleleft H_1 \triangleleft \cdots \triangleleft H_s = G \]

are two proper subnormal series for $G$. We have groups:

\[ N_{i,j} = N_i(H_j \cap N_{i+1}) \]

that form subnormal series

\[ N_i = N_{i,0} \triangleleft \cdots \triangleleft N_{i,s} = N_{i+1} \]

for all $i$, by the Butterfly Lemma. It follows that we have a subnormal series

(1) \[ 1 = N_{0,0} \triangleleft \cdots \triangleleft N_{r,s} = G.\]

Similarly, we have another subnormal series

(2) \[ 1 = H_{0,0} \triangleleft \cdots \triangleleft H_{s,r} = G \]

where

\[ H_{j,i} = H_j(N_i \cap H_{j+1}). \]

By the Butterfly Lemma, we have

\[ N_{i,j+1}/N_{i,j} \cong H_{i,j+1}/H_{i,j} \]

for all $i$ and $j$, hence the series (1) and (2) are equivalent. In particular, if a factor isomorphic to 1 occurs in one, it occurs in the other. It follows that from these we can obtain common proper refinements. \[ \square \]

If \( 1 \triangleleft N_1 \triangleleft \cdots \triangleleft N_{r-1} \triangleleft G \) is a proper subnormal series, we say the series has *length* $r$ (which equals the number of *links*). As the trivial group is not considered a simple group, any composition series for a group must be proper. Moreover, by the Correspondence Principle, no composition series can have a proper refinement.

**Theorem 17.11.** (Jordan-Hölder Theorem) Let $G$ be a group having a composition series of length $r$. Then every composition series for $G$ has length $r$ and any proper subnormal series for $G$ can be refined to a composition series for $G$. Moreover, all composition series of $G$ are equivalent.
Proof. Let \( 1 = N_0 < \cdots < N_n < G \) be a proper subnormal series for \( G \). By the Schreier Refinement Theorem, this subnormal series and a composition series for \( G \) have a common proper refinement. The result follows by the above remark. \( \square \)

Example 17.12. Let \( m = a_1 \cdots a_r \) in \( \mathbb{Z} \) with each integer \( a_i > 1 \). Then we have a subnormal series:

\[
0 = a_1 \cdots a_r \mathbb{Z}/m\mathbb{Z} < \cdots < a_1 \mathbb{Z}/m\mathbb{Z} < \mathbb{Z}/m\mathbb{Z}.
\]

As \( \mathbb{Z}/m\mathbb{Z} \) is a finite group, it has a composition series by Exercise 17.14(10) that is essentially unique. This is really the Fundamental Theorem of Arithmetic.

Remark 17.13. In an analogous way, one can define composition series for vector spaces. A “simple” vector space is a line \([\text{proof?}]\). Finite dimensional vector spaces have composition series (and infinite dimensional ones do not). Let \( V \) be a nonzero finite dimensional \( F \)-vector space. Say \( V \) has a composition series of length \( d \), i.e., we get a series \( 0 < V_1 < \cdots < V_d = V \), with simple factors \( V_i/V_{i-1} \) (after generalizing the concept of quotients to vector spaces). Thus \( d \) is the \emph{dimension} of \( V \), and the Jordan-Hölder Theorem gives an alternate proof of the invariance of the number of vectors in the basis for a finite dimensional vector space.

Exercises 17.14.
1. Prove Properties 17.4.
2. Prove that every finite abelian group is polycyclic. Use this to prove that every finite solvable group is polycyclic. Give an example of a solvable group that is not polycyclic.
3. Prove Theorem 17.3 using the isomorphism theorems. (Do not use Theorem 17.6.)
4. Prove the analogue of Theorem 17.3 replacing the word solvable with the word polycyclic using the isomorphism theorems.
5. Let \( H \) and \( K \) be subgroups of a group \( G \). Define

\[
[H, K] := \langle [h, k] \mid h \in H, k \in K \rangle.
\]

Show both of the following:

(i) If \( K \triangleleft G \) and \( K \subset H \), then \( [H, G] \subset K \) if and only if \( H/K \subset Z(G/K) \).

(ii) If \( \varphi : G \to L \) is a group epimorphism and \( A \subset Z(G) \), then \( \varphi(A) \subset Z(L) \).

6. Let \( G \) be group.
(i) Set \( \Gamma_1(G) = G \) and inductively define \( \Gamma_{n+1}(G) := [\Gamma_n(G), G] \) for \( n > 1 \). Show \( \Gamma_{n+1}(G) \subseteq \Gamma_n(G) \) and \( \Gamma_n(G) \triangleleft G \) for all \( n \).

(ii) Set \( Z^0(G) = 1 \) and inductively define \( Z^{n+1}(G) \) by

\[
Z^{n+1}(G)/Z^n(G) = Z(G/Z^n(G)),
\]

i.e., the preimage of \( G/Z^n(G) \) in \( G \) under the canonical epi-morphism. Show \( Z^{n+1}(G) \triangleleft Z^n(G) \) and \( Z^n(G) \triangleleft G \) for all \( n > 0 \).

7. Define the descending central series of \( G \) by

\[
G = \Gamma_1(G) \supset \Gamma_2(G) \supset \cdots \supset \Gamma_n(G) \supset \cdots
\]

and the ascending central series of \( G \) by

\[
1 = Z^0(G) \subset Z^1(G) \subset \cdots \subset \cdots.
\]

Prove that \( Z^n(G) = G \) if and only if \( \Gamma_n(G) = 1 \). Moreover, if this is the case, then \( \Gamma_{i+1} \subset Z^{n-i} \) for \( i = 0, \ldots, n \).

8. A group \( G \) is called nilpotent if there exists an \( n \) such that \( \Gamma_n(G) = 1 \) (and \( n \) is called the class of \( G \) if \( n \) is the least such integer). Show the following

(i) If \( G \) is nilpotent, then every subgroup of \( G \) is nilpotent.
(ii) If \( G \) is nilpotent and \( H \triangleleft G \), then \( G/H \) is nilpotent.
(iii) A direct product of finitely many nilpotent groups is nilpotent.
(iv) Give an example of a group \( G \) with \( H \triangleleft G \) with \( H \) and \( G/H \) nilpotent, but \( G \) not nilpotent.

9. Show that a group \( G \) is nilpotent if and only if there exists a normal series

\[
G = G_0 \supset G_1 \supset \cdots \supset G_n = 1
\]

in which each \( G_i \triangleleft G \) and satisfies \( G_i/G_{i+1} \subset Z(G/G_{i+1}) \) for all \( i \). Such a series is called a central series for \( G \).

10. Prove that every finite group has a composition series. Give an example of an infinite group that does not have a composition series.

11. Let \( G \) be a group and \( N \) a normal subgroup of \( G \). Prove that \( G \) has a composition series if and only if \( N \) and \( G/N \) have composition series. [It is not true that if \( G \) has a composition series that any subgroup does. Cf. Exercise 24.24(15).]

18. Free Groups

Let \( V \) be a vector space over \( F \) with basis \( \mathcal{B} \). Then \( V, \mathcal{B} \) satisfies the following universal property of vector spaces: Given any vector space \( W \) over \( F \) and diagram of set maps
with \( \text{inc} \) the inclusion of \( \mathcal{B} \) into \( V \), there exists a unique linear transformation \( T : V \to W \) such that the diagram

\[
\begin{array}{ccc}
\mathcal{B} & \xrightarrow{\text{inc}} & V \\
\downarrow f & & \downarrow T \\
W & & W
\end{array}
\]

commutes. This follows from the theorem in linear algebra that says a linear transformation \( T : V \to W \) is completely determined by where a basis for \( V \) is mapped.

Conversely, suppose that we know that \( \mathcal{B} \) is a subset of the vector space \( V \) and \( V, \mathcal{B} \) satisfies the universal property above.

Let \( W \) be the subspace spanned by \( \mathcal{B} \). The additive group \( V/W \) becomes a vector space by \( r(v + W) := rv + W \), for \( r \in F \) and \( v \in V \), with \( - : V \to W \) the natural surjective linear transformation. By the (UP) there exists a unique transformation \( T : V \to V \) satisfying

\[
\begin{array}{ccc}
\mathcal{B} & \xrightarrow{\text{inc}} & V \\
\downarrow 0 & & \downarrow T \\
V/W & &
\end{array}
\]

commutes. By uniqueness, \( T \) must be the zero map.

Let \( v \in \mathcal{B} \) and

\[
\begin{array}{ccc}
\mathcal{B} & \xrightarrow{\text{inc}} & V \\
\downarrow f_v & & \downarrow T_v \\
V & & V
\end{array}
\]

with \( f_v : \mathcal{B} \to V \) the map sending \( v \mapsto v \) and \( v' \mapsto 0 \) for all \( v' \in \mathcal{B} \{v\} \). Then there exists a unique linear transformation \( T_v : V \to V \) such that
commutes. It follows that \( B \) spans \( V \) and \( B \) is linearly independent (why?). In particular, \( V \) is completely determined by the universal property of vector spaces. Note that this says that \( V \) is completely determined by a set of generators satisfying no non-trivial relations in \( V \).

We can, in fact, modify this universal property with the inclusion map replaced by a set injection \( i_V : \mathcal{B}_0 \rightarrow V \) with \( i_V(\mathcal{B}_0) = \mathcal{B} \). Indeed, we can even drop the requirement that \( i_V \) be an injection and this still shows that \( V \) is a vector space with basis \( \mathcal{B} \). We leave this as an exercise.

We now mimic the above for groups.

**Definition 18.1.** Let \( G \) be a group with \( X \) a non-empty set in \( G \). We say that \( G \) is a free group on basis \( X \) if it satisfies the following universal property of free groups: Given any group \( G' \) and diagram

\[
\begin{array}{ccc}
X & \xleftarrow{inc} & G \\
\downarrow{f} & & \downarrow{G'} \\
\mathcal{Y} & \rightarrow & \mathcal{G'}
\end{array}
\]

of set maps with \( inc \) the inclusion of \( X \) into \( G \), then there exists a unique group homomorphism \( \varphi : V \rightarrow W \) such that the diagram

\[
\begin{array}{ccc}
X & \xleftarrow{inc} & G \\
\downarrow{f} & \downarrow{\varphi} & \\
\mathcal{Y} & \rightarrow & \mathcal{G'}
\end{array}
\]

commutes.

**Remark 18.2.** As in the vector space case, we can replace the inclusion map of \( X \) in \( G \) with a set injection \( i_G : X_0 \rightarrow G \) with \( X = i_A(X_0) \); and, in fact, the map \( i_G \) need not be assumed to be an injection.

**Remark 18.3.** Let \( X_0 \) and \( Y_0 \) be sets, \( A \) and \( B \) groups with set maps \( i_A : X_0 \rightarrow A \) and \( i_B : Y_0 \rightarrow B \) such that \( A \) is a free group on \( X = i_A(X_0) \) and \( B \) is a free group of \( Y = i_B(Y_0) \). Suppose that there exists a bijection \( j : X_0 \rightarrow Y_0 \). Then there exists unique group homomorphisms \( \varphi : A \rightarrow B \) and \( \psi : B \rightarrow A \) such that
commutes. But it then follows that $\psi \circ \varphi = 1_A$ and $\varphi \circ \psi = 1_B$, so $\varphi$ and $\psi$ are inverse isomorphisms. Thus if $A$ and $B$ are free groups on bases of the same cardinality then they are isomorphic. [Cf. If $V$ and $W$ are vector spaces over $F$ with bases of the same cardinality then they are isomorphic.]

**Theorem 18.4.** Let $X$ be a non-empty set. Then there exists a free group $G$ on $X$.

**Proof.** Choose a set $X'$ disjoint from $X$ of the same cardinality, so we have a bijection $X \to X'$ which we denote by $x \mapsto x^{-1}$. Let $X''$ be a set disjoint from $X \cup X'$ having precisely one element. Let $X'' = \{e\}$. We call $A = X \cup X' \cup X''$ the alphabet and the elements of $A$ the letters. If $x \in X$, we also write $x^1$ for $x$ and if $a$ is a letter, we also write $(a^{-1})^{-1}$ for $a$. Call a sequence $w = (a_1, \ldots, a_n, \ldots)$, $a_i \in A$, a word if $a_i \neq e$ for only finitely many $i$. Let $\tilde{W}$ be the set of words in $\times_{i=1}^{\infty} A$. The word $(e, \ldots, e, \ldots)$ is called the empty word. A word $w = (a_1, \ldots, a_n, \ldots)$ in $\tilde{W}$ is called a reduced word if $w$ satisfies

(i) $a_i \neq a_{i+1}^{-1}$ if $a_i \neq e$.
(ii) If $a_i = e$, then $a_k = e$ for all $k \geq i$.

If $w$ is a reduced word, but not the empty word, we can write it as

$$w = x_1^{n_1} \cdots x_m^{n_m}, \ x_i \in X, \ n_i \in \{\pm 1\} \text{ for all } i \text{ and some } m \text{ with } x_i^{n_i} \neq x_{i+1}^{-n_{i+1}} \text{ for all } i.$$  

We write $e$ for the empty word. Let $W$ be the set of reduced words.

**Note.** Spelling in $W$ is unique by the definition of $\times_{i=1}^{\infty} A$.

Define $\cdot : W \times W \to W$ as follows: $w = e \cdot w = w \cdot e$ for any $w \in W$ and if $w = x_1^{n_1} \cdots x_r^{n_r}$ and $w' = y_1^{m_1} \cdots y_s^{m_s}$ with $x_i, y_j \in X$ and $n_i, m_j \in \{\pm 1\}$ for all $i$ and $j$, let $w \cdot w' = x_1^{n_1} \cdots x_r^{n_r} y_1^{m_1} \cdots y_s^{m_s}$ by juxtaposition and then reduce, e.g., if $x_r^{n_r} = y_1^{-m_1}$ delete the pair unless reduced to the empty word and continue in this way to get a reduced word.

**Claim 18.5.** The map $\cdot : W \times W \to W$ is well-defined and makes $W$ into a group such that the inclusion $X \subset W$ makes $X$ into a free group on basis $X$, so $W = \langle X \rangle$. 
The details that \( \cdot : W \times W \to W \) is a well defined map is left as an exercise (cf. the above). If we show that \( W \) satisfies associativity under this map, it is then clear that \( W \) is a group. To do this we use a trick of Van der Waerden. For each \( x \in X \) and \( n \in \{ \pm \} \), define a map

\[
|x^n| : W \to W \text{ by } e \mapsto x^n \text{ and } x_1^{n_1} \cdots x_r^{n_r} \mapsto \begin{cases} x^n x_1^{n_1} \cdots x_r^{n_r} & \text{if } x^n \neq x_1^{-n_1} \\ x_2^{n_2} \cdots x_r^{n_r} & \text{otherwise.} \end{cases}
\]

Clearly,

\[
|x^n||x^{-n}| = 1_W = |x^{-n}|x^n|
\]

so \( |x^n| \) is a permutation of \( W \), i.e., \( |x^n| \in \Sigma(W) \) for all \( x \in X \) and \( n \in \{ \pm 1 \} \). Let

\[
W_0 := \langle |x| \mid x \in X \rangle \subset \Sigma(W), \text{ a subgroup.}
\]

As spelling is unique, we have a well-defined surjection

\[
\varphi : W \to W_0 \text{ given by } e \mapsto 1_W \text{ and } x_1^{n_1} \cdots x_r^{n_r} \mapsto |x_1^{n_1}| \cdots |x_r^{n_r}|.
\]

Let \( z \in W_0 \). Then there exist \( x_i \in X \) and \( n_i \in \{ \pm 1 \} \) such that

\[
z = |x_1^{n_1}| \cdots |x_r^{n_r}|. \text{ If } x_i^{n_i} = x_{i+1}^{-n_{i+1}}, \text{ then } |x_i^{n_i}| |x_{i+1}^{-n_{i+1}}| = 1_W \text{ and we can cancel them (or reduce to } 1_W), \text{ which we may assume has been done. But then}
\]

\[
\psi : W_0 \to W \text{ given by } 1_W \mapsto e \text{ and } |x_1^{n_1}| \cdots |x_r^{n_r}| \mapsto |x_1^{n_1}| \cdots |x_r^{n_r}|(e)
\]

is well-defined as \( |x_1^{n_1}| \cdots |x_r^{n_r}| = |x_1^{n_1} \cdots x_r^{n_r}| \) and spelling is unique in \( W \). In particular, \( \psi \circ \varphi = e \) and \( \varphi \circ \psi = 1_W \). It follows that \( \varphi \) is a bijection. By construction

\[
\varphi(w_1 \cdot w_2) = \varphi(w_1) \varphi(w_2) \text{ for all } c_1, w_2 \in W.
\]

As \( W_0 \) is a group, associativity holds in \( W_0 \), hence associativity holds in \( W \).

We leave it as an exercise to show that \( W \) is a free group on basis \( X \), i.e., satisfies the universal property of free groups. \( \square \)

**Exercise 18.6.** Let \( G \) be a free group on the set \( X \). Show that \( G = \langle X \rangle \).

**Definition 18.7.** Let \( H \) be a subgroup of a group \( G \). The **normal closure** of \( H \) in \( G \) is the smallest normal subgroup of \( G \) containing \( H \), i.e., the group \( \langle xHx^{-1} \mid x \in G \rangle \).
Definition 18.8. A group \( G \) is said to be defined by generators \( X = \{x_i \mid i \in I\} \) and relations \( Y = \{y_j \mid j \in J\} \) if \( G \cong E/N \) where \( E \) is a free group on basis \( X \) and \( N \) is the normal closure of \( \langle Y \rangle \). If this is the case, we say \((X,Y)\) is a presentation and also (by abuse of notation) write \( G = \langle X \mid Y \rangle \).

Corollary 18.9. Every group has a presentation, unique up to isomorphism.

**Proof.** Let \( G \) be a group and \( X \) a subset of \( G \) that generates \( G \) (e.g., we can take \( G \) itself). Let \( E \) be the free group on basis \( X \). By the universal property of free groups, there exists a unique group homomorphism \( \varphi : E \to G \) such that

\[
\begin{array}{ccc}
X & \xrightarrow{\text{inc}} & E \\
\phantom{\text{E}} & \downarrow{\varphi} & \\
\phantom{\text{X}} & G \\
\end{array}
\]

commutes. As \( G = \langle X \rangle \), the homomorphism \( \varphi \) is onto. Let \( Y \) be any subset of \( \ker \varphi \) such that \( \ker \varphi \) is the normal closure of \( \langle Y \rangle \). Then \((X,Y)\) is a presentation of \( G \). We leave the proof of uniqueness as an exercise. \( \square \)

Corollary 18.10. (Van Dyck’s Theorem) Let \( X \) be a non-empty set and \( Y \) a set of reduced words based on \( X \) (obvious definition). Let \( G = \langle X \mid Y \rangle \) be a presentation of \( G \). Suppose that \( H \) is a group satisfying \( H = \langle X \rangle \) and \( H \) satisfies all the relations in \( Y \). Then there exists a (group) epimorphism \( \psi : G \to H \).

**Proof.** Let \( E \) be the free group of \( X \) and \( N \) the normal closure of \( \langle Y \rangle \) in \( E \). As

\[
\begin{array}{ccc}
X & \xrightarrow{\text{inc}} & E \\
\phantom{\text{E}} & \downarrow{\varphi} & \\
\phantom{\text{X}} & G \\
\end{array}
\]

by the universal property of free groups, there exist unique group homomorphisms

\[
\begin{array}{ccc}
E & \xrightarrow{\theta} & G \\
\phantom{\text{E}} & \downarrow{\varphi} & \\
\phantom{\text{X}} & H \\
\end{array}
\]

We define \( \psi : G \to H \) as follows: For each \( g \in G \) choose \( x_g \) in the fiber \( \theta^{-1}(g) := \{y \in E \mid \theta(y) = g\} \) and set \( \psi(g) := \varphi(x_g) \). We must show
that \( \psi \) is well-defined, i.e., independent of the choice of \( x_g \in \theta^{-1}(g) \).
Suppose that \( x, y \) in \( E \) satisfies \( \theta(x) = g = \theta(y) \). Then \( \theta(y^{-1}x) = e_G \), so \( y^{-1}x \) lies in \( \ker \theta = N \), the normal closure of \( \langle Y \rangle \) in \( E \). If \( H = \langle X|Y' \rangle \), then \( N \) lies in the normal closure \( \ker \varphi = N' \) of \( \langle Y' \rangle \) in \( E \) by hypothesis. Consequently, \( \ker \theta \subset \ker \varphi \), so \( \varphi(x) = \varphi(y) \) and \( \psi \) if well defined. Clearly \( \psi \) is a homomorphism and surjective as both \( \theta \) and \( \varphi \) are. \( \square \)

**Examples 18.11.** Let \( Q = \langle a, b \mid a^2 = b^2, a^4 = e, bab^{-1} = a^{-1} \rangle \) be the quaternion group and \( G \) be the subgroup of \( \mathbb{M}_2(\mathbb{C}) \) generated by

\[
\begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}, \begin{pmatrix} \sqrt{-1} & 0 \\ 0 & -\sqrt{-1} \end{pmatrix}, \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}, \text{ and } \begin{pmatrix} 0 & \sqrt{-1} \\ \sqrt{-1} & 0 \end{pmatrix}.
\]

We show \( Q \cong G \). Using the Van Dyck’s Theorem, we get a group epimorphism \( Q \to G \) given by

\[
a \mapsto \begin{pmatrix} \sqrt{-1} & 0 \\ 0 & -\sqrt{-1} \end{pmatrix} \text{ and } b \mapsto \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}.
\]

Check that \( |G| \geq 8 \). The relations for \( G \) imply that \( ba = a^3b \) and \( b^3 = a^2b \), so we must have \( |Q| \leq 8 \). It follows that \( |Q| = 8 \) and \( G \cong Q \).

We end this section with a further useful example of universal properties.

**Definition 18.12.** Let \( G_1 \) and \( G_2 \) be two groups. A free product of \( G_1 \) and \( G_2 \) is a group \( E \) together with group monomorphisms

\[
(*)
\]

satisfying the following universal property: If \( H \) is a group and we group homomorphisms
then there exists a unique group homomorphism \( \varphi : E \rightarrow H \) such that the diagram

\[
\begin{array}{c}
G_1 \\
\downarrow \theta_1 \\
\psi_1 \\
\downarrow \theta_2 \\
E \\
\downarrow \varphi \\
\downarrow \psi_2 \\
G_2
\end{array}
\]

commutes.

**Theorem 18.13.** If \( G_1 \) and \( G_2 \) are groups, then a free product of \( G_1 \) and \( G_2 \) exists.

In fact if we fix the monomorphisms in (*), then the fiber product is unique up to a unique isomorphism and it is usually denoted by \( G_1 \ast G_2 \).

We leave the proof as an exercise. It is quite similar to the existence of a free group defining an alphabet in this case as follows: Assume that \( G_1 \setminus \{e_{G_1}\} \) and \( G_2 \setminus \{e_{G_2}\} \) are disjoint. Let \( \{e\} \) be a set disjoint from \( (G_1 \setminus \{e_{G_1}\}) \cup (G_2 \setminus \{e_{G_2}\}) \) and set \( A = (G_1 \setminus \{e_{G_1}\}) \cup (G_2 \setminus \{e_{G_2}\}) \cup \{e\} \).

Then define words in the obvious way, calling a non-empty word \( w \) reduced if \( w = a_1 \cdots a_r \) with \( a_i \in A \setminus \{e\} \) for \( i = 1, \ldots, r \), some \( r \) and no adjacent letters lie in the same \( G_i \), \( i = 1, 2 \). Now proceed as before.

**Exercises 18.14.**

1. Show that if a vector space \( V \) over \( F \) satisfies the universal property of vector spaces for a subset \( \mathcal{B} \), then \( \mathcal{B} \) spans \( V \) and is linearly independent. In particular, 0 cannot lie in \( \mathcal{B} \) unless \( V = 0 \).

2. Let \( V \) be a vector space over \( F \). Suppose that \( \mathcal{B}_0 \) is a set and \( i_V : \mathcal{B}_0 \rightarrow V \) a set map. Show if for every vector space \( W \) over \( F \) and set map \( f : \mathcal{B}_0 \rightarrow W \) there exists a unique linear transformation \( T : V \rightarrow W \) such that

\[
\begin{array}{c}
\mathcal{B}_0 \\
i_V \\
\downarrow f \\
V \\
\downarrow \downarrow T \\
W
\end{array}
\]

commutes, then \( i_V \) is an injective map and \( i_V(\mathcal{B}_0) \) is a basis for \( V \).

3. Prove the part of Claim 18.5 left undone.


5. Prove that a presentation for a group is unique up to isomorphism.
6. Define a free abelian group $A$ on basis $X$ to be an abelian group satisfying the following universal property of free abelian groups: Given any abelian group $B$ and diagram

$$
\begin{array}{ccc}
X & \xrightarrow{inc} & A \\
\downarrow{f} & & \downarrow{\varphi} \\
B & & B
\end{array}
$$

of set maps with $inc$ the inclusion of $X$ in $A$, there exists a unique group homomorphism $\varphi: A \to B$ such that

$$
\begin{array}{ccc}
X & \xrightarrow{inc} & A \\
\downarrow{f} & & \downarrow{\varphi} \\
B & & B
\end{array}
$$

commutes. Prove that a free abelian group on any non-empty set $X$ exists and is unique up to isomorphism.

7. Show that $\mathbb{Z}$ is a free group on basis $\{1\}$.

8. Show that if $G$ is a free group on basis $B$ and $B$ has more than two elements that $G$ is not abelian.

9. Define the free product of a set of groups $\{G_i \mid i \in I\}$ and prove it exists. It is denoted by $\ast_I G_i$.

10. Show that a free group on a set $X$ is isomorphic to $\ast_X \mathbb{Z}$.

11. Show that a free abelian group on a finite set $X$ is isomorphic to $\times_X \mathbb{Z}$. [What if $X$ is not finite?]
CHAPTER IV

Group Actions

In general, groups arise in mathematics not abstractly, but in a very concrete way. The simplest geometric example is given an object, what are all the symmetries of that object, e.g., if the object is a circle, a sphere, a tetrahedron or cube in $\mathbb{R}^3$? Of course, we also have already seen examples in the previous chapter. The dihedral group acts on a regular $n$-gon, the symmetric group acts on a set. These are all examples of a given set upon which a collection of maps acts on it bijectively. This set of maps of the object, can give us much deeper knowledge about that object. As our examples above are maps, the composition of bijections satisfies associativity and the identity map acts as the identity. This is already sufficient to give us our desired abstract formulation. In this chapter, we give many examples of this, and show how powerful an idea it is. Again we concentrate on applications to finite groups, since using the natural equivalence relation that arises allows us to count, resulting in many deep theorems about the structure of such groups.

19. The Orbit Decomposition Theorem

Groups are important because they act on objects. For example,

\begin{align*}
    \Sigma(S) & \text{ acts on } S \\
    \text{GL}_n(F) & \text{ acts on } F^n \\
    D_n & \text{ acts on the vertices of a regular } n\text{-gon,}
\end{align*}

and these actions give information about the groups and the objects on which they act. We begin by defining how a group $G$ acts on a non-empty set $S$. We then see how this leads to an equivalence relation $\sim_G$ on $S$ and form the set of equivalence classes $S/\sim_G$, followed by the application of the Mantra of Equivalence Relations. To be useful, especially when $|S|$ is finite, we must compute the cardinality of each equivalence class. Unlike the case of cosets, we cannot expect the cardinality of each equivalence class to be the same.

We begin with the definition of an action.
Definition 19.1. Let \( G \) be a group and \( S \) a non-empty set. We say that \( S \) is a (left) \( G \)-set under \( \ast : G \times S \to S \), writing \( g \ast s \) for \( \ast(g, s) \), if \( \ast \) satisfies the following: For all \( g_1, g_2 \in G \) and for all \( s \in S \),

1. \((g_1 \cdot g_2) \ast s = g_1 \ast (g_2 \ast s)\).
2. \(e \ast s = s\).

The map \( \ast : G \times S \to S \) is then called a \( G \)-action on \( S \). Note that the first property is just a generalization of associativity, the second that the unity acts like an identity.

Note. If \( S \) is a \( G \)-set and \( H \) a subgroup of \( G \), then we can also view \( S \) as an \( H \)-set by restriction, i.e., by the \( H \)-action \( \ast : H \times S \to S \) given by \( h \ast s := h \ast s \).

We now define the equivalence relation. Let \( S \) be a \( G \)-set under \( \ast \).

Define \( \sim_G \) on \( S \) by

\[
(19.2) \quad s_1 \sim_G s_2 \text{ if there exists } g \in G \text{ satisfying } s_1 = g \ast s_2.
\]

Lemma 19.3. Let \( S \) be a \( G \)-set under \( \ast \). Then \( \sim_G \) is an equivalence relation on \( S \).

Proof. Reflexitivity: For all \( s \in S \), we have \( s = e \ast s \), so \( s \sim_G s \).

Symmetry: If \( s_1 \sim_G s_2 \), then there exists a \( g \in G \) such that \( s_1 = g \ast s_2 \), hence

\[
g^{-1} \ast s_1 = g^{-1} \ast (g \ast s_2) = (g^{-1} \cdot g) \ast s_2 = e \ast s_2 = s_2.
\]

Note this shows

\[
(19.4) \quad s_1 = g \ast s_2 \text{ if and only if } g^{-1} \ast s_1 = s_2.
\]

Transitivity: Suppose that \( s_1 \sim_G s_2 \) and \( s_2 \sim_G s_3 \). Then there exist \( g, g' \in G \) satisfying \( s_1 = g \ast s_2 \) and \( s_2 = g' \ast s_3 \). Hence \( s_1 = g \ast s_2 = g \ast (g' \ast s_3) = (g \cdot g') \ast s_3 \), so \( s_1 \sim_G s_3 \). \( \square \)

Next, let \( S \) be a \( G \)-set, \( s \in S \). The equivalence class of \( s \) via \( \sim_G \) is the set

\[
\bar{s} := \{ s' \in S \mid \text{there exists } g \in G \text{ such that } s' = g \ast s \}
= \{ g \ast s \mid g \in G \}.
\]

This equivalence class is called the orbit of \( s \) under \( \ast \) and denoted by \( G \ast s \). We always let

\[
\mathcal{O} \quad \text{be a system of representatives}
\]

for the equivalence classes under \( \sim_G \).

(\( \mathcal{O} \) consists of precisely one element from each orbit — equivalence class.) We write \( G \backslash S \) for \( S/\sim_G = \{ G \ast s \mid s \in \mathcal{O} \} \), the set of orbits.
for the action $\sim_G$. Applying the Mantra of Equivalence Relations, we have

**Mantra of $G$-actions.** Let $S$ be a $G$-set, $\mathcal{O}$ a system of representatives. Then

$$S = \bigvee_{\mathcal{O}} G * s$$

and if $|S|$ is finite, then $|S| = \sum_{\mathcal{O}} |G * s|$. 

To make this useful, we must be able to compute the size of an orbit if finite. In general, orbits may have different cardinalities. We do this by attaching a group to each element of $S$ in the following way. If $S$ is a $G$-set, $s \in S$, define the stabilizer or isotropy subgroup of $s$ by

$$G_s := \{ x \in G \mid x * s = s \}.$$ 

**Lemma 19.5.** Let $S$ be a $G$-set, $s \in S$. Then $G_s$ is a subgroup of $G$.

**Proof.** By definition of a $G$-action, $e \in G_s$, so $G_s$ is non-empty. If $x, y \in G$, then $(xy) * s = x * (y * s) = x * s = s$ and if $y * s = s$, then $s = y^{-1} * s$, so $G_s$ is a subgroup. 

We can now compute the cardinality of an orbit.

**Proposition 19.6.** Let $S$ be a $G$-set, $s \in S$. Define

$$f_s : G/G_s \to G * s \text{ by } xG_s \mapsto x * s.$$ 

Then $f_s$ is a well-defined bijection. In particular, if $[G : G_s]$ is finite, then

$$|G * s| = [G : G_s] \text{ and } |G * s| \text{ divides } |G|.$$ 

**Proof.** Let $x, y \in G$. Then

$$x * s = y * s \text{ if and only if } y^{-1} * (x * s) = s$$

if and only if $(y^{-1}x) * s = s$

if and only if $y^{-1}x \in G_s$

if and only if $xG_s = yG_s$.

This shows that $f_s$ is well-defined and one-to-one. As $f_s$ is clearly surjective, the result follows.

**Example 19.7.** Let $S$ be the faces of a cube and $G$ be the group of rotations of $S$. So $G$ acts on $S$. Given any two faces $s_1, s_2$, there is an element of $G$ taking $s_1$ to $s_2$. So there is one orbit under this action. We say that $G$ acts transitively on $S$. If $s$ is a face, then the isotropy subgroup $G_s$ of $s$ is the cyclic group $C_4$ of rotations of the face $s$ about its center. By the proposition, we have

$$|G|/|C_4| = [G : G_s] = |G * s| = |S| = 6,$$
so $|G| = 24$. More generally, let $S$ be the regular solid having $n$ faces, each which has $k$ edges (or vertices) and $G$ the group of all rotations of $S$. Then the $G$-action is transitive and $|G| = nk$ by an analogous argument. It can be shown that there are only five such regular solids: the tetrahedron ($n = 4, k = 3$), the cube, the octahedron ($n = 8, k = 3$), the dodecahedron ($n = 12, k = 5$), and the icosahedron ($n = 20, k = 3$). So the corresponding rotation groups have 12, 24, 24, 60, 60 elements and are isomorphic to $A_4$, $S_4$, $S_4$, $A_5$, $A_5$ respectively. We give further details in Section 20.

Suppose that $S$ is a $G$-set, $s \in S$. We say that $G \ast s$ is a one point orbit of $S$ and $s$ is a fixed point (under the action of $G$) if $G \ast s = \{s\}$. We set

$$F_G(S) := \{s \in S \mid |G \ast s| = 1\} \subset S,$$

the set of fixed points of $S$ under the action of $G$.

If $S$ is a $G$-set, a major problem is to determine if $F_G(S)$ is non-empty, and if it is, to compute $|F_G(S)|$.

The following equivalent conditions characterize the one point orbits. The proof is left as an easy exercise.

**Lemma 19.8.** Let $S$ be a $G$-set, $s \in S$. Then the following are equivalent:

1. $s \in F_G(S)$.
2. $G_s = G$.
3. $G \ast s = \{s\}$.

In particular, if $\mathcal{O}$ is a system of representative, then $F_G(S) \subset \mathcal{O}$.

If $\mathcal{O}$ is a system of representatives for a $G$-action on a set $S$, we always set

$$\mathcal{O}^* = \mathcal{O} \setminus F_G(S).$$

In particular, if $s \in \mathcal{O}^*$, we have $|G \ast s| = [G : G_s] > 1$, i.e., $G_s < G$. (Recall that the symbol $<$ when used for sets means a subset, but not the whole set.)

Putting this all together, we get the desired result:

**Theorem 19.9.** (Orbit Decomposition Theorem) Let $S$ be a $G$-set. Then

$$S = F_G(S) \lor \bigvee_{\mathcal{O}^*} G \ast s.$$

In particular, if $S$ is a finite set, then

$$|S| = \sum_{\mathcal{O}} |G \ast s| = |F_G(S)| + \sum_{\mathcal{O}^*} [G : G_s].$$

[Note if $s \in \mathcal{O}^*$, then $1 < [G : G_s]$ and if $G$ is finite then $[G : G_s] \mid |G|$.]
20. Addendum: Finite Rotation Groups in \( \mathbb{R}^3 \).

In this addendum, we give further details to Example 19.7. Recall if \( R \) is a commutative ring, then the orthogonal group is \( \text{O}_n(R) := \{ A \in \text{GL}_n(R) \mid AA^t = I \} \) and the special orthogonal group is \( \text{SO}_n(R) = \text{O}_n(R) \cap \text{SL}_n(R) \). In this section, we shall determine the finite subgroups of \( \text{SO}_3(\mathbb{R}) \). Elements of \( \text{SO}_3(\mathbb{R}) \) are very special as any element is, in fact, a rotation about some axis through the origin in \( \mathbb{R}^3 \). This is easy to see using linear algebra. Indeed let \( A \) be a nonidentity element of \( \text{SO}_3(\mathbb{R}) \). Viewing \( \mathbb{R}^3 \) as an inner product space under the dot product \( \cdot \) the matrix \( A \) is an isometry, i.e., \( A \) preserves dot products as \( Av \cdot Aw = v \cdot A^tAw = v \cdot w \) for all \( v \) and \( w \) in \( \mathbb{R}^3 \) and also in \( \mathbb{C}^3 \). It follows that every eigenvalue \( \varepsilon \) of \( A \) in \( \mathbb{C} \) satisfies \( |\varepsilon| = 1 \). In particular, \( \varepsilon \) is a root of unity, so is \( \pm 1 \) if real. As the characteristic polynomial \( f_A \) of \( A \) is of odd degree, it has a real root, hence an eigenvalue 1 or \(-1\). If all the roots of \( f_A \) are real, det \( A = 1 \) forces \( A \) to be similar to \( \text{diag}(1, -1, -1) \), a rotation of angle \( \pi \) around some axis. If the roots of \( f_A \) are not all real, they must be \( 1, e^{2\pi \sqrt{-1}/n}, e^{-2\pi \sqrt{-1}/n} \), some \( n \), and \( A \) is similar to the rotation \[
\begin{pmatrix}
1 & 0 & 0 \\
0 & \cos 2\pi/n & \sin 2\pi/n \\
0 & -\sin 2\pi/n & \cos 2\pi/n
\end{pmatrix}
\] about some axis.

We view \( \text{SO}_3(\mathbb{R}) \) acting on the unit sphere \( S^2 \). So if \( A \) is a nonidentity element in \( \text{SO}_3(\mathbb{R}) \), it is the rotation by some angle about some axis through the origin. This axis intersects \( S^2 \) in two points, that we call the poles of \( A \). If \( A \) is of finite order, then the angle of rotation must be \( 2\pi/n \) for some integer \( n \). Let \( G \) be a finite subgroup of \( \text{SO}_3(\mathbb{R}) \) of order \( n \), and set \( P := \{ p \mid p \text{ is a pole of some element } A \text{ in } G \} \), called the set of poles of \( G \).
Lemma 20.1. Let $G$ be a finite subgroup of $\SO_3(\mathbb{R})$ of order $n$, and $P$ the set of poles of $G$. Then $P$ is a $G$-set under the natural action. In particular, if $p \in P$, $|G| = |G : G_p||G_p|$ and $|G * p| = |G : G_p|$, where $G_p$ is the isotropy subgroup of $p$ and $G * p$ the orbit of $p$.

**Proof.** Let $p$ be a pole of $A$ in $G$. If $B$ is a matrix in $G$, then $(BAB^{-1})Bp = Bp$, so $Bp$ is a pole of $BAB^{-1}$. □

Let $G$ be a finite subgroup of $\SO_3(\mathbb{R})$ of order $n$ and $P$ the set of poles. If $p$ is a pole, let $n_p = |G_p|$. As every nonidentity element of $G$ has two poles and each nonidentity element in $G_p$ has $p$ as a pole, we have

$$2(n - 1) = \sum_P (n_p - 1).$$

Let $G * p_1, \ldots, G * p_r$ be the (distinct) orbits of $P$ and $n_i = n_{p_i}$, for $i = 1, \ldots, r$. As the order of stabilizers of elements in a fixed orbit are all the same, we have

$$2n - 2 = \sum_{i=1}^{r} (n_i - 1)|G * p_i|,$$

so dividing by $n = |G|$ implies that

$$2 - \frac{2}{n} = \sum_{i=1}^{r} (n_i - 1) \frac{|G * p_i|}{n} = \sum_{i=1}^{r} (n_i - 1) \frac{1}{|G_{p_i}|} = \sum_{i=1}^{r} (1 - \frac{1}{n_i}).$$

As $n > 1$ and $1 - (1/n_i) \geq 1/2$, we have $1 \leq 2 - (2/n) \leq 2$. It follows that we can only have $r = 1, 2, 3$, i.e., there are at most three orbits. We investigate each possibility separately.

**Case.** $r = 1$:

This cannot occur, as $1 \leq 2 = (2/n) = 1 - (1/n_1)$ is impossible.

**Case.** $r = 2$:

As

$$2 - \frac{2}{n} = (1 - \frac{1}{n_1}) + (1 - \frac{1}{n_2}),$$

we have

$$\frac{2}{n} = \frac{1}{n_1} + \frac{1}{n_2}.$$

It follows that $n = n_1 = n_2$ and $|G * p_1| = |G * p_2| = 1$ as $n = n_i|G * p_i|$. In particular, $|P| = 2$. Since every nonzero element in $G$ has two poles, we must have $G$ is a cyclic group of finite order with axis connecting the two poles.

**Case.** $r = 3$:
In this case, we have

\[ \frac{2}{n} = \frac{1}{n_1} + \frac{1}{n_2} + \frac{1}{n_3}, \]

and we may assume that \( n_1 \leq n_2 \leq n_3 \). We must find the groups with data \((n_1, n_2, n_3)\). Let \([|G * p_1|, |G * p_2|, |G * p_3|] = [n/n_1, n/n_2, n/n_3]\) be the tuple of the sizes of the corresponding orbits. If \( n_1 \geq 3 \), then \( 2/n \leq 0 \), which is impossible, so \( n_1 = 2 \), and

\[ \frac{1}{2} + \frac{2}{n} = \frac{1}{n_2} + \frac{1}{n_3}. \]

This is impossible if \( n_2 \geq 4 \), so \( n_2 = 2 \) or \( 3 \).

**Subcase.** \( r = 3, n_2 = 2 \).

By (*) we must have \( n_3 = n/2 \). Therefore, \( G_{p_3} \) has two poles so is cyclic of order \( n/2 \). As \( P \) is a \( G \)-set and the nonzero element \( g \) in \( G_{p_1} \) must be a rotation of angle \( \pi \), it follows that the poles of \( g \) must be perpendicular to the poles of the nonzero elements in \( G * p_3 \) (and reverse pole pairs). It follows that \( G \) must be the dihedral group \( D_\frac{n}{2} \).

Therefore we may assume that \( n_2 = 3 \). If \( n_3 \geq 6 \), then \( (1/2) + (1/3) + (1/n_3) \leq 0 \), so \( n_3 = 3, 4, \) or \( 5 \).

**Subcase.** \( G \) has data \((2, 3, 3)\):

It follows by (*) that \( n = 12 \), so \( G \) also has orbit data \([6, 4, 4]\). As \( P \) is a \( G \)-group, if \( p \) lies in \( G * p_3 \) and \( q \) is the closest point to \( p \) in \( G * p_2 \), then \( G_{p_3}q \) consists of at least three points equidistance from \( p \), which form an equilateral triangle. It follows that there are four such triangles with the poles forming a regular tetrahedron. Therefore, we see that \( G \) is the rotation group of a regular tetrahedron and the orbit data \([6, 4, 4]\) is the number of edges, vertices, and faces of the tetrahedron, respectively. (Cf. Example 19.7.) This group can be shown to be isomorphic to \( A_4 \) using Exercise 24.24(7) below.

**Subcase.** \( G \) has data \((2, 3, 3)\):

It follows by (*) that \( n = 24 \), so \( G \) also has orbit data \([12, 8, 6]\). In this case, if \( p \) lies in \( G * p_3 \) and \( q \) is the closest point to \( p \) in \( G * p_2 \), then \( G_{p_3}q \) consists of at least four points equidistance from \( p \), which form a square, and we see that \( G \) is the group of rotations of a cube with orbit data \([12, 8, 6]\), the number of edges, vertices, and faces of a cube, respectively. (Cf. Example 19.7.) This group can be shown to be isomorphic to \( S_4 \) using Exercise 24.24(7) below.

**Subcase.** \( G \) has data \((2, 3, 5)\):

It follows by (*) that \( n = 60 \), so \( G \) also has orbit data \([30, 20, 12]\). In this case, if \( p \) lies in \( G * p_3 \) and \( q \) is the closest point to \( p \) in \( G * p_2 \),
then $G_{p,q}$ consists of at least five points equidistance from $p$, which form a pentagon, and we see that $G$ is the group of rotations of a regular icosahedron with orbit data $[30, 20, 12]$, the number of edges, vertices, and faces of a regular icosahedron, respectively. (Cf. Example 19.7.) This group can be shown to be isomorphic to $A_5$ using Exercise 24.24(7) below.

21. Examples of Group Actions

In this section, to demonstrate the usefulness of this concept, we give many examples of group actions. As this is the most important concept in our study of groups, there are many exercises at the end of this section. To understand this topic, one should do many of them. In the following two sections, we shall give further specific applications of some of examples given in this section.


In this example, we let the $G$-set $S$ be $G$ itself. The (left) action is given by

$*: G \times S \to S$ by $g * s = gsg^{-1}$

called conjugation by $G$. The orbit of an element $a \in S = G$ is

$C(a) := G * a = \{xax^{-1} \mid x \in G\}$,

called the conjugacy class of $a$, and the stabilizer of $a$ is

$Z_G(a) := G_a = \{x \in G \mid xax^{-1} = a\} = \{x \in G \mid xa = ax\}$.

This subgroup of $G$ is called the centralizer of $a$. So $Z_G(a)$ is the set of elements of $G$ commuting with $a$. In particular, $\langle a \rangle \subset Z_G(a)$. The set of fixed points of this action is

$(21.2) \quad F_G(S) = \{a \in S \mid xax^{-1} = a \text{ for all } x \in G\}$

$= \{a \in G \mid xa = ax \text{ for all } x \in G\}$,

the center $Z(G)$ of $G$. Recall it is a normal subgroup of $G$. Of course, in general, a $G$-set is not a group, so one cannot expect the set of fixed points to have an algebraic structure.

Let $C$ be a system of representatives for the conjugation action of $G$ on $G$. So $C^* = C \setminus Z(G)$, and the Mantra of $G$-actions becomes

$G = Z(G) \lor \bigvee_{C^*} C(a),$

with

$(21.3) \quad |G| = |Z(G)| + \sum_{C^*} [G : Z_G(a)],$
if \( G \) is finite. Equation 21.3 is called the class equation. We give an interesting application of the class equation. If \( G \) is a finite group of order \( p^n \), for some (positive) prime \( p \) and \( n > 0 \), we call \( G \) a \( p \)-group. So a subgroup of a \( p \)-group is either a \( p \)-group or the trivial group.

**Application 21.4.** If \( G \) is a \( p \)-group, then \( 1 < Z(G) \subseteq G \).
(Recall this means \( 1 \neq Z(G) \).) In particular, if \( G \) is a \( p \)-group of order \( p^n \), with \( n > 1 \), then \( G \) is not simple:

We know if \( a \in C^* \), then \( p \mid [G : Z_G(a)] \), so

\[
0 \equiv |G| = |Z(G)| + \sum_{a \in C^*} [G : Z_G(a)] \equiv |Z(G)| \mod p.
\]

Since \( e \in Z(G) \), we must have \( |Z(G)| \geq 1 \), so \( p \mid |Z(G)| \) and \( |Z(G)| \geq p \). In particular, \( 1 < |Z(G)| \subset G \). If \( G \) is not abelian, then \( 1 < Z(G) < G \), and \( G \) is not simple, as \( Z(G) \triangleleft G \). If \( G \) is abelian, then \( G = Z(G) \) and there exists an element \( a \) in \( G \) of order \( p \) by Exercise 10.16(2). As \( |G| > p \), we have \( 1 < \langle a \rangle < G \) with \( \langle a \rangle \triangleleft G \).

**Useful Observation.** If we are trying to prove a property about groups and \( G \) is a finite non-abelian group with \( Z(G) > 1 \) (which in general may not occur), then \( G/Z(G) \) is a group of order less than that of \( G \) and the canonical epimorphism \( \pi : G \to G/Z(G) \) may allow us to apply induction.

**Computation 21.5.** Let \( G = D_3 = \{e, r, r^2, f, fr, rf\} \), with \( |G| = 6 \) and satisfying \( r^3 = e = f^2 \) and \( frf^{-1} = r^{-1} = r^2 \). We have \( fr = r^2f \) and \( rf = fr^2 \), so

\[
C(e) = \{e\} \quad \text{and} \quad 1 \mid 6
\]
\[
C(r) = \{r, r^2\} \quad \text{and} \quad 2 \mid 6
\]
\[
C(f) = \{f, fr, rf\} \quad \text{and} \quad 3 \mid 6
\]

and

\[
Z_G(e) = G \quad \text{and} \quad |C(e)| = [G : Z_G(e)] = 1
\]
\[
Z_G(r) = \{e, r, r^2\} \quad \text{and} \quad |C(r)| = [G : Z_G(r)] = 2
\]
\[
Z_G(f) = \{e, f\} \quad \text{and} \quad |C(f)| = [G : Z_G(r)] = 3
\]

with fixed points

\[
Z(G) = \{e\},
\]

so

\[
|G| = |Z(G)| + |C(r)| + |C(f)| = 1 + 2 + 3.
\]
Note that the conjugacy classes (equivalence classes), $C(r) = C(r^2)$ and $C(f) = C(fr) = C rf$. This application leads to the following theorem.

**Theorem 21.6.** Let $G$ be a group of order $p^2$ with $p$ a prime. Then $G$ is abelian.

**Proof.** Suppose that $G$ is not abelian, then $Z(G) < G$ by Exercise 21.26(1). By Application 21.4, we have $1 < Z(G)$; so by Lagrange’s Theorem, we must have $|Z(G)| = p$. Let $a \in G \setminus Z(G)$. Then we have $Z(G) < Z_G(a)$, as $a \in Z_G(a)$. But then the subgroup $Z_G(a)$ must satisfy $|Z_G(a)| = p^2 = |G|$ by Lagrange’s Theorem, hence $Z_G(a) = G$. This implies that $a \in Z(G)$, a contradiction. \[\square\]

**Remarks 21.7.** Let $G$ be a $p$-group.

1. If $G = p^2$, the group $G$ is abelian, but may not be cyclic. If $G$ is cyclic, then $G \cong \mathbb{Z}/p^2\mathbb{Z}$. However, the group $\mathbb{Z}/p\mathbb{Z} \times \mathbb{Z}/p\mathbb{Z}$ is not cyclic.

2. $D_4$ is a 2-group of order $2^3$, but not abelian.

Later we shall prove a major theorem. It will imply that every finitely generated abelian group, which includes all finite abelian groups, is a direct product of cyclic groups. There is also a uniqueness statement. At present, we only can show that abelian groups of order $pq$, with $p, q$ (not necessarily distinct) primes are of this form.

We now jazz up Example 21.1 by letting $G$ act on a subgroup $H$ by conjugation. For this to work, however, we must have $xHx^{-1} \subseteq H$ for all $x \in G$, i.e., we must have $H \triangleleft G$. So suppose that $H \triangleleft G$. Let $G$ act on $H$ by conjugation, i.e.,

$$*: G \times H \rightarrow H \text{ is given by } g * h = ghg^{-1}.$$  

We still have $C(a) = G * a$ for all $a \in H$, so if $H$ is a finite group, we have

$$|H| = |F_G(H)| + \sum_{a \in \mathcal{O}} [G : Z_G(a)],$$

as $Z_G(a) = G_a := \{x \in G \mid xax^{-1} = a\}$, with $a \in \mathcal{O}$, where $\mathcal{O}$ is a system of representatives for the $G$-action on $H$. An application of this is:

**Proposition 21.8.** Let $G$ be a $p$-group and $N$ a non-trivial normal subgroup of $G$. Then there exists an non-identity element $x \in N$ such that $xy = yx$ for all $y \in G$, i.e., $1 < N \cap Z(G)$. [Cf. the case $N = G$ with Application 21.4.]
Proof. Let $G$ act on $N$ by conjugation. Every subgroup of $G$ is either 1 or a $p$-group, hence
\[ 0 \equiv |N| = |F_G(N)| + \sum_{a \in \mathcal{O}} [G:Z_G(a)] \equiv |F_G(N)| \mod p \]
by the Orbit Decomposition Theorem. As $xx^{-1} = e$ for all $x \in G$, the unity $e$ is a fixed point, so $p \mid |F_G(N)| \geq 1$. Therefore, $|F_G(N)| \geq p$. The result now follows from
\[ F_G(N) = \{ z \in N \mid xzx^{-1} = z \ \text{for all} \ x \in G \} = \{ z \in N \mid xz = zx \ \text{for all} \ x \in G \} = N \cap Z(G). \]
We can even generalize this conjugation type action a bit further as follows: Let $H$ be a subgroup of $G$. Define the normalizer of $H$ in $G$ by
\[ N_G(H) := \{ x \in G \mid xHx^{-1} = H \}. \]
The normalizer has the following properties:

Properties 21.9. Let $H$ be a subgroup of $G$. Then the following are true:
1. $N_G(H)$ is a subgroup of $G$.
2. $H \trianglelefteq N_G(H)$.
3. If $H \trianglelefteq K$ with $K$ a subgroup of $G$, then $K \subset N_G(H)$.
4. $N_G(H)$ is the unique largest subgroup of $G$ containing $H$ as a normal subgroup.

Replacing $G$ by $N_G(H)$, we can let $N_G(H)$ act on $H$ by conjugation, by the above, i.e., we have an $N_G(H)$-action
\[ * : N_G(H) \times H \to H \ \text{given by} \ x \ast h = xhx^{-1}. \]
In fact, we can go one step further. Let $K$ be a subgroup of $G$ and set $A = K \cap N_G(H)$. Then the action
\[ * : A \times H \to H \ \text{given by} \ x \ast h = xhx^{-1} \]
makes $H$ into an $A$-set. Instead of pursuing this, we generalize conjugation in another direction.

Example 21.10. Conjugation on Sets.

Let $S = \mathcal{P}(G) := \{ A \mid A \subset G, \ \text{a subset} \}$, the power set of $G$. Then $S$ becomes a $G$-set by
\[ * : G \times S \to S \ \text{given by} \ x \ast A = xAx^{-1} \]
where \( xA x^{-1} = \{ xax^{-1} \mid a \in A \} \). Note that \( |A| = |xA x^{-1}| \) for all \( x \in G \). We also call this action \textit{conjugation}. The orbit of \( A \in P(G) \) is the \textit{conjugacy class}

\[
C(A) = G \ast A = \{ xA x^{-1} \mid x \in G \},
\]

and the stabilizer of \( A \) is

\[
N_G(A) := G_A = \{ x \in G \mid xA x^{-1} = A \},
\]
called the \textit{normalizer} of \( A \) in \( G \). The Orbit Decomposition Theorem yields

\[
S = P(G) = \bigvee_{C} C(A) = F_G(S) \vee \bigvee_{C} C(A).
\]

If \( G \) is finite so is \( P(G) \), and we would then have

\[
2^{|G|} = |S| = |F_G(S)| + \sum_{C} [G : N_G(A)].
\]

This is not so useful, so we try to cut \( S \) down. We use the following:

\textbf{Observation 21.11.} Let \( S \) be a non-empty set and \( \sim \) an equivalence relation of \( S \). Suppose a subset \( \emptyset \neq T \subset S \) satisfies the following:

For all \( a \in T \), we have \([a]_{\sim} \subset T\).

Then \( \sim \mid_T \) (actually \( \sim \mid_{\tau \times \tau} \)) defines an equivalence relation on \( T \) and

\[
T / \sim_T = \{ [a]_{\sim} \mid a \in T \}.
\]

[We are just saying: If \( C \) partitions \( S \) and \( \emptyset \neq D \subset C \), then \( D \) partitions \( \bigcup_{D} \overline{a} \).]

There are many interesting subsets of \( P(G) \). Recall that if \( H \) is a subgroup of \( G \), then so is \( xHx^{-1} \) for all \( x \in G \).

\textbf{Specific Examples 21.12.} Let \( G \) be a group. Then the following subsets of \( P(G) \) are \textit{G}-sets by conjugation (i.e., restricting the action above) if non-empty:

1. \( P(G) \).
2. \( \mathcal{G} := \{ H \mid H \subset G \text{ a subgroup } \} \).
3. \( \mathcal{T}_n := \{ A \mid A \subset G \text{ with } |A| = n \} \).
4. \( \mathcal{G} \cap \mathcal{T}_n \).
5. If \( G \) is a finite group, \( p \) a (positive) prime satisfying \( p^r || |G| \), with \( r > 0 \), so \( G = p^r m \) with \( p \) and \( m \) relatively prime, then

\[
\text{Syl}_p(G) := \mathcal{G} \cap \mathcal{T}_{p^r} = \{ H \mid H \text{ a subgroup of } G \text{ of order } p^r \}.
\]
6. If \( H \in \mathcal{G} \) then \( C(H) \).
7. If \( A \in P(G) \) then \( C(A) \).
We can also push this further by restricting the action of $G$ on one of these sets to the action by a subgroup. We shall later study this action in greater detail, especially (5) and (6).

**Example 21.13. Translation Action.**

Let $G$ be a group, then the power set $\mathcal{P}(G)$ is a $G$-set via

\[ * : G \times \mathcal{P}(G) \to \mathcal{P}(G) \]  
\[ \text{defined by } g * T = gT := \{ gx \mid x \in T \} \]

called *translation*. As before, this is more interesting when we replace $\mathcal{P}(G)$ by appropriate subsets. For example, using the notation of (21.12), $G$ acts on $T_n$ by translation.

**Warning 21.14.** $G$ does not act on $G$ by translation as $H \in G$ does not mean that $x * H \in G$. In fact, $x * H$ would lie in $G$ if and only if $x \in H$.

One also can translate other sets by $G$. For example, if $H$ is a subgroup of $G$, we can let $G$ act on the cosets $G/H$ by translation, i.e., let

\[ * : G \times G/H \to G/H \]  
\[ \text{be given by } x * aH = xaH. \]

We have used this action before, viz., in the proof of the General Cayley Theorem. The stabilizer of $aH$ in $G/H$ under translation by $G$ is given by:

\[ G_{aH} = \{ x \in G \mid xaH = aH \} = \{ x \in G \mid a^{-1}xaH = H \} \]
\[ = \{ x \in G \mid a^{-1}xa \in H \} = \{ x \in G \mid x \in aHa^{-1} \} \]
\[ = aHa^{-1}. \]

Note that if $H < G$ in the above then $F_G(G/H) = \emptyset$ as $G = G_{aH}$ would imply that $G = aHa^{-1}$.

The following proposition is a useful computing devise. We leave its proof as an exercise.

**Proposition 21.15.** Let $G$ be a finite group, $H$ a subgroup of $G$. Suppose that $H$ is a $p$-group and $p \mid [G : H]$. Then $p \mid [N_G(H) : H]$. In particular, $H < N_G(H)$. Consequently, if $G$ is a finite $p$-group and $H < G$ a subgroup, then $H < N_G(H)$.

We finish the discussion of translation actions to answer our outstanding question on how cosets arise naturally. Indeed, if $H$ is a subgroup of $G$, and $G$ is a right $H$-set by the right $H$-action by translation on $G$, i.e.

\[ * : G \times H \to G \]  
\[ \text{given by } g * h = gh, \]
then the orbit $g * H$ is the left coset $gH$ and the Orbit Decomposition Theorem in this case is just Lagrange’s Theorem.

**Example 21.16. Endomorphic Action.**

Let $V$ be a vector space over a field $F$. A linear transformation $T : V \to V$ is called an endomorphism of $V$. The set

$$\text{End}_F(V) := \{ T \mid T : V \to V \text{ a linear transformation} \}$$

is a ring under the usual addition and composition of functions. [From linear algebra, you should know the relationship between $\text{End}_F(V)$ and $\mathbb{M}_n(F)$ if $V$ is n-dimensional.] Let $G$ be $\text{End}_F(V)$ viewed as an additive group. Then we have $V$ a $G$-set by

$$*: G \times V \to V \text{ given by } T * v = T(v)$$
called evaluation. If $V = F^n$, we can replace $\text{End}_F(V)$ by $\mathbb{M}_n(F)$.

Much more interesting is the next example.

**Example 21.17. Automorphic Action.**

Let $V$ be a vector space over a field $F$ and $G = \text{Aut}_F(V) = (\text{End}_F(V))^\times$. Then $V$ is a $G$-set via the evaluation

$$*: G \times V \to V \text{ given by } T * v = T(v).$$

**Example 21.18. Addendum to Example 21.1.**

Let $R$ be any ring, $S = \mathbb{M}_n(R)$, and $G = \text{GL}_n(R)$. Then $S$ is a $G$-set under conjugation, i.e.,

$$* : G \times S \to S \text{ is given by } A * B = ABA^{-1}.$$ 

This is an important action and leads to the following:

**Problem 21.19.** Let $F$ be a field and $*$ the action above. Find a nice system of representatives for this action.

If every non-constant polynomial with coefficients in the field $F$ has a root in $F$, we say that $F$ is algebraically closed, e.g., the complex numbers are algebraically closed by the Fundamental Theorem of Algebra (to be proven later). The above sought after system of representatives is called the set of Jordan canonical forms. For a general field $F$, the sought after set is the set of rational canonical forms. Proving this will be a major goal later.

[We have another action by

$$* : G \times S \to S \text{ given by } A * B = ABA^t.$$]
and ask the same questions. This is especially interesting if \( S = M_n(F) \) is replaced by the subset of symmetric matrices or the subset of skew symmetric matrices.

**Example 21.20. Evaluation.**

We have previously seen another example of an action given by evaluation, viz., if \( S \) is a non-empty set, then \( S \) is a \( \Sigma(S) \)-set by the evaluation map

\[
\ast : \Sigma(S) \times S \to S \text{ given by } \gamma \ast s = \gamma(s).
\]

If \( H \) is a subgroup, we can restrict the action to \( H \). For example, let \( \gamma \in \Sigma(S) \) and \( \Gamma = \langle \gamma \rangle \). Then \( S \) becomes a \( \Gamma \)-set by the action\[\ast : \Gamma \times S \to S \text{ given by the evaluation } \gamma^i \ast s = \gamma^i(s) \text{ for all integers } i.\]

We shall use this action later.

It should be mentioned that evaluation actions are used all the time. If you have an object, then its automorphism group acts on it by evaluation. For example, such objects can be algebraic as above, topological, or geometric.

**Example 21.21. Pullback Action.**

As mentioned before a \( G \)-action on a set \( S \) can be restricted to the action of any subgroup of \( G \). In the examples above, we have used this. We now generalize this. Recall that a subgroup of a group is just a group in which the inclusion map is a monomorphism. More generally, suppose that we have a group homomorphism \( \theta : G \to G' \) and a \( G' \)-set \( S \) defined by \( \ast : G' \times S \to S \). Then \( S \) becomes a \( G \)-set by the pullback action defined by

\[
\ast : G \times S \to S \text{ is given by } g \ast s = \theta(g) \ast s.
\]

**Example 21.22. Shift Action.**

Our last example, is an action that leads to a nice result called Cauchy’s Theorem, which is the first step toward showing that finite groups have subgroups of the order of the power of a prime for any power of that prime dividing the order of the group. Let \( G \) be a finite group and \( p \) a (positive) prime dividing the order of \( G \). Set

\[
S = \{(g_1, \ldots, g_p) \in G^p \mid g_1 \cdots g_p = e \text{ in } G\} \subset G^p,
\]

where \( G^p = G \times \cdots \times G \) (\( p \) times). As \( (e, \ldots, e) \in S \), the set \( S \) is non-empty. Moreover, it follows immediately that

\[
(g_1, \ldots, g_p) \in S \implies g_p = (g_1 \cdots g_{p-1})^{-1} \text{ in } G.
\]
Therefore, for all ordered \((p - 1)\)-tuples \(g_1, \ldots, g_{p-1}\), with \(g_i \in G\), we have a unique \(g_p\) such that \((g_1, \ldots, g_p)\) lies in \(S\). Therefore, \(|S| = |G^{p-1}|\), so \(p \mid |S|\). If \(y \in \mathbb{Z}\), let \(\tilde{y} \in [y]_p = \tilde{y}\) denote the smallest positive integer in \(\tilde{y}\). With this notation, if \(x \in \mathbb{Z}\), the ordered tuple

\[
\tilde{1} + x, \ldots, \tilde{p} + x
\]

will be a (cyclic) permutation of \(1, \ldots, p\). Let \((\mathbb{Z}/p\mathbb{Z}, +)\) act on \(S\) by

\[
\pi \ast (g_1, \ldots, g_p) = (g_{1+x}, \ldots, g_{p+x})
\]

i.e., if \(1 \leq i, x \leq p\), then the action of \(\pi\) shifts \(i\) by \(x\) slots to the number \(x + i\) until it hits \(p\) which it then maps to 1, etc. For example, if \(p = 5\), then \(\tilde{3} \ast (g_1, g_2, g_3, g_4, g_5) = (g_4, g_5, g_1, g_2, g_3)\). The right hand side of (*) still lies in \(S\), since \(g_1 \cdots g_p = e\) implies \(g_{i+1} \cdots g_p g_1 \cdot g_i = e\) by conjugating successively by \(g_i^{-1}, \ldots, g_i^{-1} \ (1 \leq i < p)\). It follows easily that \(S\) is a \(\mathbb{Z}/p\mathbb{Z}\)-set under this action. With this action, we can now prove:

**Theorem 21.23.** (Cauchy’s Theorem) *Let \(p\) be a (positive) prime dividing the order of the finite group \(G\). Then there exists an element of \(G\) of order \(p\).*

**Proof.** Let \(S\) be the \(\mathbb{Z}/p\mathbb{Z}\)-group under the shift action of Example 21.22. Apply the Orbit Decomposition Theorem 19.9 to get

\[
|G^{p-1}| = |S| = |F_{\mathbb{Z}/p\mathbb{Z}}(S)| + \sum_{\mathcal{O}^*} [\mathbb{Z}/p\mathbb{Z} : (\mathbb{Z}/p\mathbb{Z})_x].
\]

If \(x \in \mathcal{O}^*\), then \((\mathbb{Z}/p\mathbb{Z})_x = 1\) as \(\mathbb{Z}/p\mathbb{Z}\) has only the trivial subgroups. So

\[
0 \equiv |S| \equiv |F_{\mathbb{Z}/p\mathbb{Z}}(S)| \mod p.
\]

As \((e, \ldots, e) \in F_{\mathbb{Z}/p\mathbb{Z}}(S)\), we have \(|F_{\mathbb{Z}/p\mathbb{Z}}(S)| \geq p > 1\). Thus there exists an element \((e, \ldots, e) \neq (g_1, \ldots, g_p) \in F_{\mathbb{Z}/p\mathbb{Z}}(S)\). But this means that

\[
(g_1, \ldots, g_p) = (g_2, \ldots, g_p, g_1) = \cdots = (g_{i+1}, \ldots, g_p, \ldots, g_i) = \cdots = (g_p, g_1, \ldots, g_{p-1})
\]

lies in \(G^p\). It follows that \(e \neq g_1 = \cdots = g_p\). Consequently, if \(g = g_1\), then \((g, \ldots, g) \in S\), which means that \(g^p = e\) as needed. \(\Box\)

**Corollary 21.24.** *Let \(p\) be a prime and \(G\) a non-trivial finite group. If every non-identity element in \(G\) has order a power of \(p\), then \(G\) is a \(p\)-group.*

**Corollary 21.25.** *Let \(p\) and \(q\) be primes and \(G\) of order \(pq\). Then \(G\) is not simple.*
21. EXAMPLES OF GROUP ACTIONS

Proof. If \( p = q \), then \( G \) is abelian hence not simple as \( G \) is not of prime order. If \( p < q \), then \( G \) contains a normal group of order \( q \) by Useful Counting 12.7. □


1. Let \( G \) be a group. Show all of the following:
   (a) \( Z(G) = \bigcap_{a \in G} Z_G(a) \).
   (b) \( a \in Z(G) \) if and only if \( C(a) = \{a\} \) if and only if \( |C(a)| = 1 \).
   (c) \( a \in Z(G) \) if and only if \( \mathbb{G} = Z_G(a) \).
   (d) If \( G \) is finite, then \( a \in Z(G) \) if and only if \( |Z_G(a)| = |G| \).

2. Let \( G \) be a group and \( k(G) \) the number of conjugacy classes in \( G \).
   Suppose that \( G \) is finite. Show that \( k(G) = 3 \) if and only if \( G \) is isomorphic to the cyclic group of order three or the symmetric group on three letters.

3. Let \( G \) be a group of order \( p^2 \) with \( p \) a prime. Show that either \( G \cong \mathbb{Z}/p^2\mathbb{Z} \) or \( G \cong \mathbb{Z}/p\mathbb{Z} \times \mathbb{Z}/p\mathbb{Z} \).

4. Show that Properties 21.9 are valid.

5. Compute all the conjugacy classes and isotropy subgroups in \( A_4 \).

6. Let \( G \) be a group of order \( p^n \), \( p \) a prime. Suppose the center of \( G \) has order at least \( p^{n-1} \). Show that \( G \) is abelian.

7. Let \( G \) be a \( p \)-group. Using Exercise 17.14(8), show that \( G \) is a nilpotent group. In particular a finite product of \( p \)-groups for various \( p \) is nilpotent.

8. Let \( G \) be a nilpotent group (cf. 17.14(8) and \( H \) a proper subgroup of \( G \). Show that \( H \neq N_G(H) \).

9. Suppose that \( H \) is a proper subgroup of a finite group \( G \). Show that \( G \neq \bigcup_{g \in G} gHg^{-1} \). [This may not be true if \( G \) is infinite.]

10. Let \( H \) be a subgroup of \( G \). Let \( H \) act on \( G/H \) by translation. Compute the orbits, stabilizers, and fixed points of this action.

11. Prove Proposition 21.15. [Hint use the previous exercise.]

12. Let \( G \) be a group and \( S \) a non-empty set. Show
   (a) If \( \ast : G \times S \rightarrow S \) is a \( G \)-action, then \( \varphi : G \rightarrow \Sigma(S) \) given by \( \varphi(x)(s) = x \ast s \) (i.e., if we let \( \varphi_x = \varphi(x) \) then \( \varphi_x(s) = x \ast s \) is a group homomorphism (called the permutation representation).
   (b) If \( \varphi : G \rightarrow \Sigma(S) \) is a group homomorphism then \( \ast : G \times S \rightarrow S \) given by \( x \ast s = \varphi_x(s) \) where \( \varphi_x = \varphi(x) \) is a \( G \)-action.

13. Let \( Q \) be the quaternion group. Show that there exist no group monomorphism \( \varphi : Q \rightarrow S_7 \).
14. Let $G$ be a group of order $pq$ with $q \leq p$ primes. Show that $G$ is not simple and has a normal subgroup of order $p$.

15. Let $G$ be a finite $p$-group. Show if $p^n \mid |G|$ then $G$ has a normal subgroup of order $p^n$.

16. Let $p$ be an odd prime. Classify all groups $G$ of order $2p$ up to isomorphism.

17. Let $S$ be a $G$-set. Suppose that both $S$ and $G$ are finite. If $x \in G$ define the fixed point set of $x$ in $G$ by

$$F_S(x) := \{ s \in S \mid x \cdot s = s \}.$$ 

Show the number of orbits $N$ of this action satisfies

$$N = \frac{1}{|G|} \sum_G |F_S(x)|.$$ 

[So $N$ is the average of the size of the fixed point sets of elements of $G$.]

22. Sylow Theorems

Let $G$ be a finite group of order of order $p^r m$ with $p$ a prime relatively prime to $m$ and $r > 0$. A subgroup $H$ of $G$ is called a Sylow $p$-subgroup of $G$ if $H$ has order $p^r$. Let

$$\text{Syl}_p(G) := \{ H \mid H \text{ a Sylow } p\text{-subgroup of } G \}.$$ 

The natural question is answered by the next result. Its proof uses much of what we have done.

**Theorem 22.1.** (First Sylow Theorem) Let $p$ be a (positive) prime dividing the order of a finite group $G$. Then $\text{Syl}_p(G)$ is non-empty, i.e., there exists a Sylow $p$-subgroup.

We prove the stronger

**Theorem 22.2.** (Generalized First Sylow Theorem) Let $p$ be a (positive) prime such that $p^s$, with $s \geq 1$, divides the order of a finite group $G$. Then there exists a subgroup of $G$ of order $p^s$.

**Proof.** We induct on the order of $G$. If $|G| \leq p$, the result is trivial. We make the following:

**Induction Hypothesis.** If $T$ is any finite group satisfying $|T| < |G|$ and $p^s \mid |T|, s \geq 1$, then $T$ contains a subgroup of order $p^s$.

We must show that $G$ contains a subgroup of order $p^s$. If $H < G$ is a subgroup of $G$ satisfying $p \nmid [G : H]$, then $p^s \mid |H|$. By the Induction
Hypothesis, $H$, hence $G$, contains a subgroup of order $p^s$. Therefore, we are done unless we make the following:

**Assumption.** If $1 < H < G$ is a subgroup, then $p \mid [G : H]$.

For example, if $a \in G$ satisfies $Z_G(a) < G$, then $p \mid [G : Z_G(a)]$. In particular, this is true for any $a \notin Z(G)$, i.e., $a \in C^* = C \setminus Z(G)$, where $C$ is a system of representatives for conjugate action of $G$ on $G$ in Example 21.1. Applying the Class Equation 21.3, we have

$$0 \equiv |G| = |Z(G)| + \sum_c [G : Z_G(a)] \equiv |Z(G)| \mod p,$$

so $p \mid |Z(G)|$. By Cauchy’s Theorem 21.23, there exists an element $a \in Z(G)$ of order $p$. Let $H = \langle a \rangle$, a subgroup of $G$ of order $p$. Since $a \in Z(G)$, we have $xa^ix^{-1} = a^i$ for all $x \in G$ and $i \in \mathbb{Z}$. In particular, $H \triangleleft G$ (cf. Example 11.5(3)) and $G/H$ is a group of order $[G : H] = |G|/|H| = |G|/p < |G|$. Consequently, $p^{s-1} \mid |G/H|$; and, by the Induction Hypothesis, there exists a subgroup $T \subset G/H$ of order $p^{s-1}$. Let $\tilde{T} : G \to G/H$ be the canonical epimorphism. Then $\tilde{T}$ is the kernel and by the Correspondence Principle, there exists a subgroup $\tilde{T}$ of $G$ containing $H$ and satisfying $\tilde{T} = T/H$. Hence $|\tilde{T}| = |T||H| = p^s$, and $\tilde{T}$ works.

To investigate the set of Sylow $p$-subgroups, $\text{Syl}_p(G)$ of $G$, we need three lemmas.

**Lemma 22.3.** Let $H$ and $K$ be subgroups of a finite group $G$. Then $|HK| = |H||K|/|H \cap K|$.

**Proof.** This is Equation 13.5. (Cf. Exercise 10.16(4).) □

The main mathematical reason for the validity of the further Sylow theorems below is the following:

**Lemma 22.4.** Let $G$ be a finite group, $P$ a Sylow $p$-subgroup, and $H$ a subgroup of $G$ that is also a $p$-group. If $H$ is a subgroup of $N_G(P)$, i.e., $hPh^{-1} = P$ for all $h \in H$, then $H$ is a subgroup of $P$. In particular, if, in addition, $H$ is also a Sylow $p$-subgroup of $G$, then $H = P$.

**Proof.** The key to the proof of this lemma is the Second Isomorphism Theorem together with the counting version of the Second Isomorphism Theorem of Lemma 22.3. By definition, $P \triangleleft N_G(P)$ and by hypothesis $H \subset N_G(P)$. Hence, by the Second Isomorphism Theorem, we have $HP$ is a subgroup of $N_G(P)$, with $P \triangleleft HP$ and $H \cap P \triangleleft H$, and satisfying $HP/P \cong H/H \cap P$. In particular, by Lemma 22.3, we...
have
\[(*) \quad |HP| = |P| \frac{|H|}{|H \cap P|}.\]
As \(P, H\) are \(p\)-groups and \(H \cap P \subset H\) a subgroup, (*) implies that \(HP\) is a \(p\)-group satisfying \(P \subset HP \subset G\). Thus \(HP\) is a Sylow \(p\)-subgroup of \(G\), i.e., \(|P| = |HP|\), so \(P = HP\), hence \(H \subset P\). If, in addition, \(|H| = |P|\) then \(H = P\). □

We convert the lemma above into the language of group actions with the action given by conjugation on sets. Explicitly, we use Specific Examples 21.12(6) after restricting the action of \(G\) to a subgroup.

**Lemma 22.5.** Let \(P\) be a Sylow \(p\)-subgroup of a finite group \(G\) and \(H\) a subgroup of \(G\) that is also a \(p\)-group. Then the following hold:

1. \(C(P)\) consists of Sylow \(p\)-subgroups of \(G\).
2. The conjugacy class \(C(P)\) is an \(H\)-set by conjugation, i.e., via the action \(*: H \times C(P) \to C(P)\) by conjugation.
3. Suppose that \(T\) is a fixed point under the \(H\)-action in (2), i.e.,
   \[T \in F_H(C(P)) = \{W \in C(P) \mid xWx^{-1} = W \quad \text{for all} \ x \in H\}.\]
   Then \(H \subset T\).
4. If \(H\) in (3) is a Sylow \(p\)-subgroup, then, under the \(H\)-action in (2), \(H\) is the only possible fixed point, i.e., \(F_H(C(P)) \subset \{H\}\).

**Proof.**

1. As \(|P| = |xPx^{-1}|\) for all \(x \in G\), this is clear.
2. is immediate.
3. As \(T\) is an \(H\)-fixed point, the orbit \(H \ast T = \{T\}\), so \(xT^{-1} = T\) for all \(x \in H\). By definition, this means that \(H \subset N_G(T)\). By Lemma 22.4 and (1), we conclude that \(H \subset T\).
4. If \(T \in F_H(C(P))\), then \(H \subset T\) by (3). If further, \(H\) is a Sylow \(p\)-subgroup, then by (3), we have \(H = T\). □

**Note.** Under the hypothesis of Lemma 22.5, we have the isotropy subgroup of \(P\) is given by \(H_P = N_G(P) \cap H\) for any Sylow \(p\)-subgroup of \(G\). But \(P\) is then also a Sylow \(p\)-subgroup of \(N_G(P)\), so by restricting the \(H\)-action to \(H_P\) the proof of Lemma 22.5 shows that \(N_G(P) \cap H \subset P\).

**Theorem 22.6.** (Sylow Theorems). Let \(G\) be a finite group, \(p\) a (positive) prime dividing the order of \(G\). Then the following are true:

**First Sylow Theorem.** There exists a Sylow \(p\)-subgroup of \(G\).

**Second Sylow Theorem.** All Sylow \(p\)-subgroups are conjugate, i.e., if \(P\) is a Sylow \(p\)-subgroup of \(G\), then \(C(P) = \text{Syl}_p(G)\).
**Third Sylow Theorem.** Let $P$ be a Sylow $p$-subgroup of $G$. Then the following are true:

(i) $|\text{Syl}_p(G)| = [G : N_G(P)]$.

(ii) $|\text{Syl}_p(G)| = |G|$.

(iii) $|\text{Syl}_p(G)| \equiv 1 \mod p$.

(iv) $|\text{Syl}_p(G)| = |G : P|$, i.e., if $|G| = p^n m$ with $p \nmid m$, then $|\text{Syl}_p(G)| \mid m$.

**Fourth Theorem.** Let $H$ be a subgroup of $G$ that is a $p$-group. Then $H$ lies in some Sylow $p$-subgroup of $G$.

**Proof.** Let $H$ be a subgroup of $G$ that is a $p$-group. We have already proven the First Sylow Theorem, so let $P$ be a Sylow $p$-subgroup of $G$. By Lemma 22.5, we have $C(P) \subset \text{Syl}_p(G)$ is an $H$-set under conjugation. By the Orbit Decomposition Theorem, we have

$$|C(P)| = |F_H(C(P))| + \sum_{O^*} [H : H_T].$$

If $T \in O^*$, then $H_T < H$, so $p \mid [H : H_T]$, hence

$$|C(P)| \equiv |F_H(C(P))| \mod p.$$  

**Note.** Equation (†) is interesting, since $H$ does not appear at all on the left hand side, i.e., (†) is true for all $p$-subgroups of $H$. To evaluate a constant function (of such $p$-subgroups of $H$) given by (†), it suffices to evaluate at any such $H$.

**Case 1.** Let $H = P$:

By Lemma 22.5, we conclude that $F_P(C(P)) \subset \{P\}$. As $xPx^{-1} = P$ for all $x \in P$, we must have $F_P(C(P)) = \{P\}$.

We use Case 1, to evaluate the constant in (†) as

$$|C(P)| \equiv |F_P(C(P))| = 1 \mod p.$$  

**Case 2.** Let $H$ be a Sylow $p$-subgroup (not necessarily $P$):

Applying (†) and (*), we have

$$1 \equiv |C(P)| \equiv |F_H(C(P))| \mod p.$$  

Hence $F_H(C(P))$ is non-empty so $F_H(C(P)) = \{H\}$ by Lemma 22.5. In particular, $H \in F_H(C(P)) \subset C(P)$ proving that $\text{Syl}_p(G) = C(P)$, which is the Second Sylow Theorem. Plugging this information into (*), we have

$$|\text{Syl}_p(G)| = |C(P)| \equiv 1 \mod p.$$  

Since $|C(P)| = [G : N_G(P)] = |G|/|N_G(P)|$ and $P \subset N_G(P)$ with $p \nmid [G : N_G(P)]$, we also have proven the Third Sylow Theorem.
Case 3. Let $H$ be an arbitrary $p$-subgroup of $G$: By (*), we have

$$|F_H(C(P))| \equiv |C(P)| \equiv 1 \mod p.$$ 

Thus there exists a $T \in F_H(C(P))$, hence $H \subseteq T \in \text{Syl}_p(G)$ by Lemma 22.5. This proves the Fourth Theorem.

Recall that a subgroup $H$ of a group $G$ is called a characteristic subgroup of $G$ if the restriction map $\text{res} : \text{Aut}(G) \to \text{Aut}(H)$ given by $\varphi \mapsto \varphi|_H$ is well-defined, i.e., the target is correct. We write $H \trianglelefteq \triangleleft G$. We know that characteristic subgroups are normal, and being characteristic is a transitive property.

**Proposition 22.7.** Let $G$ be a finite group, $p$ a (positive) prime dividing the order of $G$, and $P$ a Sylow $p$-subgroup. Then the following are equivalent:

1. $\text{Syl}_p(G) = \{P\}$.
2. $C(P) = \{P\}$.
3. $P \triangleleft G$.
4. $P \triangleleft \triangleleft G$.

**Proof.** The equivalence of (1), (2), and (3) follows from the Second Sylow Theorem. Clearly, (1) implies (4) as all isomorphic images of a group have the same cardinality. Finally (4) implies (3) by the comment above.

**Corollary 22.8.** Let $G$ be a finite group, $p$ a (positive) prime dividing the order of $G$, and $P$ a Sylow $p$-subgroup. Then $N_G(P) = N_G(N_G(P))$.

**Proof.** As $P$ is a Sylow $p$-subgroup of $N_G(P)$ and normal in it, $P$ is a characteristic subgroup of $N_G(P)$ by the proposition. As $N_G(P)$ is a normal subgroup of $N_G(N_G(P))$, we conclude by Exercise 11.9(20) that $P$ is normal in $N_G(N_G(P))$, i.e., $N_G(N_G(P)) \subseteq N_G(P)$. So we must have equality.

A stronger result is left as Exercise 22.15(1) below.

**Examples 22.9.** Let $G$ be a finite group, $p$ a prime dividing the order of $G$. In the following examples, we shall use the following notation:

$P_p$ will denote an arbitrary Sylow $p$-subgroup of $G$.

$n_p$ will denote the number of Sylow $p$-subgroups.
So there exists an integer \( k \) (depending on \( p \)) with

\[
n_p = |\text{Syl}_p(G)| = [G : N_G(P)] = 1 + pk
\]

\( P \) is normal if and only if \( k = 0 \).

1. Every group of order 15 is cyclic (hence isomorphic to \( \mathbb{Z}/15\mathbb{Z} \)) and every group of order 45 is abelian (hence isomorphic to \( \mathbb{Z}/45\mathbb{Z} \) or \( \mathbb{Z}/15\mathbb{Z} \times \mathbb{Z}/3\mathbb{Z} \)).

In both cases, the Third Sylow Theorem implies the Sylow \( p \)-subgroups are normal for \( p = 3, 5 \). As groups of order \( p \) are cyclic and of order \( p^2 \) are abelian, the result follows by Exercise 13.7(9). [We leave the statement about isomorphism as an exercise.]

2. If \( G \) is a group of order \( p^rq \) with \( p > q \) primes, then \( P \) is normal by the Third Sylow Theorem.

3. If \( G \) is a group of order \( 315 = 3^2 \cdot 5 \cdot 7 \), then \( G \) is not simple:

Suppose this is false. Then no Sylow \( p \)-subgroup is normal. We have

\[
n_3 = 1 + 3k \mid 5 \cdot 7 \text{ some } k, \text{ so } [G : N_G(P_3)] = 7.
\]

\[
n_5 = 1 + 5k \mid 3^2 \cdot 7 \text{ some } k, \text{ so } [G : N_G(P_5)] = 21.
\]

By Lagrange’s Theorem,

\[
N_G(P_3) \text{ is of order 45 and } N_G(P_5) \text{ is of order 15,}
\]

so both are abelian, and \( \text{Syl}_5(N_G(P_3)) \subset \text{Syl}_5(G) \). In particular, if \( T \in \text{Syl}_5(N_G(P_3)) \), we have \( T < N_G(P_3) \) (or \( 1 + 5k \mid 3^2 \) implies \( k = 0 \)). Thus \( N_G(P_3) \subset N_G(T) \). We conclude that

\[
45 = |N_G(P_3)| \mid |N_G(T)| = 15,
\]

which is impossible.

4. If \( G \) is a group of order \( 525 = 3 \cdot 5^2 \cdot 7 \), then \( G \) is not simple:

Suppose this is false. Then no Sylow \( p \)-subgroup is normal. Let \( P \) and \( Q \) be distinct Sylow 5-subgroups and set \( I = P \cap Q \). We must have \( |I| = 1 \) or 5. Although \( PQ \) is not a group, we know by Exercise 10.16(4) that

\[
|PQ| = |P| |Q| / |P \cap Q| = |P| |Q| / |I|.
\]

Suppose that \( I = 1 \). Then \( |PQ| = 5^4 > |G| \), so this is impossible. Hence \( |I| = 5 \). Both \( P \) and \( Q \) are abelian as of order \( 5^2 \), so \( P, Q \subset N_G(I) \). Therefore,

\[
|N_G(I)| \geq |PQ| = 5^3.
\]

As \( |N_G(I)| \mid |G| \), it follows that

\[
|N_G(I)| = 3 \cdot 5^2 \cdot 7 \text{ or } 5^2 \cdot 7.
\]
In the first case, $I \triangleleft N_G(I) = G$; and in the second case $[G : N_G(I)] = 3$ and $|G| \not| 3!$, so $N_G(I)$ contains a non-trivial normal subgroup of $G$ by Useful Counting 12.7. In fact, $N_G(I) \triangleleft G$, as 3 is the smallest prime dividing $|G|$ (cf. Corollary 12.10).

5. If $G$ is a group of order $945 = 3^3 \cdot 5 \cdot 7$, then $G$ is not simple:

Suppose this is false. Then no Sylow $p$-subgroup is normal. As $n_3 = 1 + 3k \mid 5 \cdot 7$, we have $[G : N_G(P_3)] = n_3 = 7$. Since $|G| \not| 7!$, the subgroup $N_G(P_3)$ contains a non-trivial normal subgroup of $G$ by Useful Counting 12.7.

6. If $G$ is a group of order $1785 = 3 \cdot 5 \cdot 7 \cdot 17$, then $G$ is not simple:

Suppose this is false. Then no Sylow $p$-subgroup is normal. We have

\[
n_{17} = 1 + 17k \mid 3 \cdot 5 \cdot 7, \quad n_{17} = 5 \cdot 7.
\]
\[
n_3 = 1 + 3k \mid 5 \cdot 7 \cdot 17, \quad n_3 = 7, 5 \cdot 17, \text{ or } 5 \cdot 7 \cdot 17.
\]

As $|G| = 1785 \not| 7!$, we can eliminate $n_3 = 7$ by Useful Counting 12.7. Thus

\[
|N_G(P_{17})| = 3 \cdot 17 \quad \text{and} \quad |N_G(P_3)| = 3 \cdot 7 \quad \text{or} \quad 3.
\]

As $3 \mid |N_G(P_{17})|$ and $3 \mid |G|$, we may assume that $P_3 \subset N_G(P_{17})$. But $1 + 3k \mid 17$ implies $k = 0$, so $P_3 \triangleleft N_G(P_{17})$. Consequently, $N_G(P_{17}) \subset N_G(P_3)$, implying that

\[
3 \cdot 17 = |N_G(P_{17})| \leq |N_G(P_3)| = 3 \cdot 7 \quad \text{or} \quad 3,
\]

which is impossible.

7. Here is a fact useful in counting. Suppose that $|G| = pm$ with $p \not| m$, $p$ a prime, then

$G$ contains $n_p(p-1)$ elements of order $p$, so $|G| \geq n_p(p-1)+1$, adding in the identity. Similarly, if $|G| = pqm$ with $p \not| m$, $q \not| m$, and $p$ and $q$ distinct primes, then

\[
|G| \geq n_p(p-1) + n_q(q-1)+1,
\]

adding in the identity. This last remark does not work if $p = q$, because two Sylow $p$-groups can intersect non-trivially, and it is difficult to estimate the number of elements in the union of all the Sylow $p$-subgroups.

For some other results, see the exercises below.

The proof of the following theorem illustrates many of the ideas that we have used in the examples.

**Theorem 22.10.** Let $p$ and $q$ be primes, $G$ a group of order $p^s q$ with $s \geq 1$. Then $G$ is not simple.
PROOF. We may assume that $p \neq q$ and there exist no normal Sylow subgroups. In particular, if $n_p = |\text{Syl}_p(G)|$, then $1 < n_p = 1 + kp \mid q$ with $k \geq 1$, an integer. In particular, $n_p \leq q$. Among all the Sylow $p$-subgroups of $G$, choose two distinct ones, $P_1$ and $P_2$, satisfying $I = P_1 \cap P_2$ has maximal order. If $I = 1$, then all Sylow $p$-subgroups of $G$ intersect pairwise in the identity, so the number of non-identity elements in the Sylow $p$-subgroups is $(p^s - 1)q$. As $|G| = p^s q$, we conclude that the Sylow $q$-group is normal, a contradiction. So we may assume that $1 < I$. As $I$ is a $p$-group but not a Sylow $p$-subgroup, we have $1 < I < N_{P_i}(I)$ for $i = 1, 2$ by Proposition 21.15.

Let $H = \langle N_{P_1}(I), N_{P_2}(I) \rangle$, the group generated by $N_{P_1}(I)$ and $N_{P_2}(I)$. Then we must have $I \triangleleft H$.

**Claim.** $q \mid |H|$. Suppose not. Then $H$ is a $p$-group, so by the Fourth Theorem, there exists a Sylow $p$-group $P_3$ of $G$ such that $H \subset P_3$. Hence for $i = 1, 2$, we have

$$P_1 \cap P_2 = I < N_{P_i}(I) \subset P_i \cap H \subset P_i \cap P_3.$$ 

This contradicts the maximality of $|I|$, lest $P_1 = P_1 \cap P_3 = P_2 \cap P_3 = P_2$. This proves the claim, i.e., $q \mid |H|$. Let $Q$ be a Sylow $q$-group of $H$, hence also of $G$ as $q \mid |G|$. Consequently, by Exercise 10.16(4), $|P_1 Q| = |P_1| |Q| / |P_1 \cap Q| = |P_1| |Q| = |G|$ as $p$ and $q$ are relatively prime. This means that

$$(*) \quad G = P_1 Q := \{xy \mid x \in P_1, y \in Q\}.$$ 

If $g \in G$, set $I^g := gIg^{-1}$ and $N = \langle I^x \mid x \in G \rangle$. By Great Trick 11.3, we have $I \subset N \triangleleft G$. Let $g \in G$. By $(*)$, we can write $g = xy$, for some $x \in P_1$ and $y \in Q$. In particular, as $y \in Q \subset H \subset N_G(I)$, we have

$$I^g = I^{xy} = xyIy^{-1}x^{-1} = xIx^{-1} = I^x \subset P_1,$$

as $I \subset P_1$. Thus $N \subset P_1 < G$, so $1 < N \triangleleft G$, and the theorem follows. \hfill \Box

**Definition 22.11.** Let $G$ be a group. We say that a subgroup $M$ of $G$ is a **maximal subgroup** if $M < G$ and if $M < H \subset G$ is a subgroup, then $H = G$. We say that $G$ is called an **(internal) direct product** of normal subgroups $N_1, \ldots, N_r$ if $G = N_1 \cdots N_r$ and $N_i \cap N_1 \cdots N_{i-1} N_{i+1} \cdots N_r = 1$, for $i = 1, \ldots, r$. (Cf. Exercises 13.7(3)(4)(5).)

An interesting class of groups in the finite case is determined by the next result theorem.

**Theorem 22.12.** Let $G$ be a finite group. Then the following are equivalent
(1) Every Sylow subgroup of $G$ is normal in $G$.
(2) $G$ is isomorphic to an (internal) direct product of its Sylow subgroups.
(3) Every maximal subgroup of $G$ is normal in $G$.

**Proof.** Let $|G| = p_1^{e_1} \cdots p_r^{e_r}$ be a factorization of $|G|$ with distinct primes $p_1, \ldots, p_r$ and $e_i > 0$, $i = 1, \ldots, r$. Let $P_i \in \text{Syl}_{p_i}(G)$ for $i = 1, \ldots, r$. So $P_i$ is normal in $G$ if and only if $\text{Syl}_{p_i}(G) = \{P_i\}$.

(1) $\Rightarrow$ (2): Since $P_1 \cdots P_n$ is a subgroup of $G$ of the same order as $G$ they are equal. This is an internal direct product as $p_i$ and $p_1 \cdots p_{i-1}p_{i+1} \cdots p_r$ are relatively prime and all the $P_i$ are normal subgroups of $G$.

(2) $\Rightarrow$ (3): Let $M$ be a maximal subgroup of $G$. By Exercise 13.7 (6), $M$ is a direct product of the $M \cap P_i$, $i = 1, \ldots, r$. Since $M$ is maximal, it must be $P_1 \cdots P_{i-1}P_{i+1} \cdots P_r$ for some $i$.

(3) $\Rightarrow$ (1): If $P_i$ is not normal in $G$, then $P < N_G(P_i) < G$ by Corollary 22.8. As $G$ is a finite group, there exists a maximal subgroup $M$ of $G$ containing $N_G(P_i)$ and hence also $P_i \in \text{Syl}_{p_i}(M)$. As $P_i \subset N_G(P_i) \subset M$ and $M \vartriangleleft G$, by the Frattini Argument (Exercise 22.15(2) below), we have $G = N_G(P_i)M$. But $N_G(P_i)M \subset M$, a contradiction. 

**Definition 22.13.** A finite group $G$ is called nilpotent if it satisfies the equivalent conditions of Theorem 22.12 are satisfied. This agrees with the definition of a nilpotent group defined in Exercise 17.14(nilpotent-groups) for finite groups.

**Corollary 22.14.** Every finite nilpotent group is solvable.

**Exercises 22.15.**
1. Let $G$ be a finite group, $p$ a (positive) prime dividing the order of $G$, and $P$ a Sylow $p$-subgroup. Suppose that $H$ is a subgroup of $G$ satisfying $N_G(P) \subset H \subset G$, then $H = N_G(H)$.
2. (Frattini Argument) Let $G$ be a finite group and $K \vartriangleleft G$. If $P$ is a $p$-Sylow subgroup of $K$, show that $G = KN_G(P)$.
3. Let $G$ be a finite group of order $p^2q$ with $p$ and $q$ primes. Show that $G$ is not simple without using Theorem 22.10.
4. Let $G$ be a finite abelian group of order $p^2q$ with $p$ and $q$ distinct primes. Show that $G$ is isomorphic to either $\mathbb{Z}/p^2q\mathbb{Z}$ or $\mathbb{Z}/p\mathbb{Z} \times \mathbb{Z}/pq\mathbb{Z}$ and the second group is isomorphic to $\mathbb{Z}/p\mathbb{Z} \times \mathbb{Z}/p\mathbb{Z} \times \mathbb{Z}/q\mathbb{Z}$.
5. Let $G$ be a finite group of order $p^2q^2$ with $p$ and $q$ primes. Show that $G$ is not simple.
6. Let $G$ be a group of order 56. Using Example 22.9(7) but not Theorem 22.10, show that $G$ is not simple.

7. Let $G$ be a finite group of order $pqr$ with $p$, $q$, and $r$ primes. Show that $G$ is not simple.

8. Let $G$ be a finite group, $p$ the smallest prime dividing the order of $G$. Suppose that $G$ contains a cyclic Sylow $p$-subgroup. Show that $N_G(P) = Z_G(P)$. (Cf. Exercise 12.11(7)).

9. Let $p$ be a prime dividing the order of a finite group $G$. Let $P \in \text{Syl}_p(G)$. If $N \triangleleft G$ show that $N \cap P \in \text{Syl}_p(N)$.

10. Let $G$ be a finite group with $p$ a prime dividing $|G|$ but not a $p$-group. Suppose that $G$ contains two Sylow $p$-subgroups whose intersection $I$ has maximal cardinality. Show the $N_G(I)$ is not a $p$-group.

11. Let $G$ be a finite group with $p$ a prime dividing $G$ but not a $p$-group. Suppose that $G$ contains two Sylow $p$-subgroups whose intersection is not trivial. Show that there exists a subgroup $H < G$ such that either $G = N_G(H)$ or $H$ is not a $p$-subgroup but contains a Sylow $p$-subgroup of $G$.

12. Show that no group of order 144 is simple.

13. Show that no group of order 2000 is simple.

14. Using Exercise 17.14(8) and Exercise 21.26(7) show that a finite group is nilpotent if and only if it is a product of its Sylow subgroups.

15. A normal subgroup $N \neq 1$ of $G$ is called a minimal normal subgroup if it contains no normal subgroups of $G$ other than itself and 1. Show if $G$ is a nontrivial finite solvable group, then any minimal normal subgroup of $G$ is an elementary $p$-group, i.e., a group isomorphic to $(\mathbb{Z}/p\mathbb{Z})^n$ for some prime $p$ and positive integer $n$.

16. If $G$ is a finite nilpotent group, then the center of $G$ is not trivial. In particular, $S_3$ is not nilpotent.

17. If $G$ is a finite nilpotent group, then every subgroup and quotient group of $G$ is nilpotent. Show the converse is false in general.

18. Let $G$ be a finite group. If every proper subgroup of $G$ is nilpotent, then $G$ is solvable.

23. Addendum: Finite Solvable Groups

In this section, we wish to extend the Sylow theorems in the case of a finite solvable group. We use this to give a new characterization of finite solvable groups. The proofs give a nice application of results that we have proven.

We will be interested in the following subgroups.
**Definition 23.1.** Let $G$ be a finite group of order $mn$ with $m$ and $n$ relatively prime. A subgroup of $G$ of order $n$ is called an *Hall $n$-subgroup* of $G$.

Of course, every Sylow $p$-subgroup is a Hall subgroup. We are interested when Hall $n$-subgroups exist and if any two such are conjugate. In general, this is false. We shall prove that this is true in the case that $G$ is a solvable group.

**Definition 23.2.** Let $G$ be a non-trivial group. A subgroup $1 < N \subset G$ is called a *minimal normal subgroup* of $G$ if there exists no normal subgroup $H$ of $G$ with $1 < H < N$.

Note that if $G$ is a non-trivial finite group, then minimal normal subgroups of $G$ exist. Of course, if $G$ is simple, then $G$ is the minimal normal subgroup of $G$.

In the following proposition, we use properties about the internal direct product of finitely many groups. [Cf. Exercise 13.7(3)-(6).]

**Proposition 23.3.** Let $G$ be a finite group and $1 < N \subset G$ be a minimal normal subgroup. Among all the subgroups of $G$ isomorphic to $N$, choose those $N = N_1, \ldots, N_r, r \geq 1$ that satisfy

$$H = N_1 \cdots N_r \cong N_1 \times \cdots \times N_r$$

and with $|H|$ maximal.

Then $H$ is a characteristic subgroup of $G$.

**Proof.** We know that $H \triangleleft G$. Let $\sigma : G \to G$ be a group automorphism. We have to show that $\sigma(N_i) \subset H$ for $i = 1, \ldots, r$. We know that $\sigma(N_i)$ is a normal subgroup of $G$ isomorphic to $N$, so if $\sigma(N_i)$ is not a subgroup of $H$, then $\sigma(N_i) \cap H < \sigma(N_i)$, hence is of lesser order. Since $\sigma(N_i) \cap H$ is a normal subgroup of $G$, the minimality of $|H|$ implies that $\sigma(N_i) \cap H = 1$. But this implies that the subgroup generated by $H$ and $\sigma(N_i)$ is isomorphic to $H \times \sigma(N_i)$ (why?), contradicting the maximality of $|H|$. Therefore, $\sigma(N_i) \subset H$, and $H$ is a characteristic subgroup of $G$. \qed

**Corollary 23.4.** Let $G$ be a non-trivial finite group having no non-trivial characteristic subgroups. Then $G$ is either a simple group or a direct product of isomorphic simple groups.

**Proof.** We use the notation of Proposition 23.3. If there exists a normal subgroup $1 < H \triangleleft N$, then $H \triangleleft N_1 \cdots N_r$, since $N_1N_2 \cdots N_r \cong N_1 \times \cdots \times N_r$. By Proposition 23.3 and the hypothesis, we must have
As \( N \) is a minimal normal subgroup of \( G \), we must have \( N = H \), so \( N \) is simple. The result follows. \( \square \)

**Corollary 23.5.** Let \( N \) be a minimal normal subgroup of finite group \( G \). Then \( N \) is either simple group or a direct product of isomorphic simple groups.

**Proof.** If \( 1 < H \) is a characteristic subgroup of \( N \), then \( N \) is a normal subgroup of \( G \) by Exercise 11.9(20). The result follows by Corollary 23.4. \( \square \)

**Corollary 23.6.** Let \( N \) be a minimal normal subgroup of a finite solvable group \( G \). Then there exists a prime \( p > 0 \) such that \( N \) is cyclic of order \( p \) or the direct product of cyclic groups of order \( p \). In particular, \( N \) is abelian.

**Proof.** This follows immediately as groups of prime order constitute the abelian simple groups. \( \square \)

We need a simple lemma whose proof we leave as an exercise.

**Lemma 23.7.** Let \( G \) be a group and \( H \subset G \) a subgroup. Then for all \( x \in G \), we have \( \text{N}_G(xHx^{-1}) = x\text{N}_G(H)x^{-1} \).

We state and prove the first key result in this section. It generalizes the Sylow theorems in the case of a finite solvable group.

**Theorem 23.8.** (Hall) Let \( G \) be a finite solvable group of order \( mn \) with \( m \) and \( n \) relatively prime. Then

1. There exists a Hall \( m \)-subgroup of \( G \).
2. Every subgroup of \( G \) whose order divides \( m \) lies in a Hall \( m \)-subgroup of \( G \).
3. Every pair of Hall \( m \)-subgroups of \( G \) are conjugate.
4. Suppose that \( |G| = p_1^{e_1} \cdots p_r^{e_r} \) is a factorization into distinct primes with \( e_i > 0, i = 1, \ldots, n \). Then the number \( h_m \) of Hall \( m \)-subgroups of \( G \) is a product of factors each of which is congruent to \( 1 \) modulo some prime factor of \( m \) and a power of a prime dividing some \( p_i^{e_i}, i = 1, \ldots, r \).

**Proof.** As subgroups and homomorphic images of solvable groups are solvable by Theorem 17.3, this allows us to induct on the order of a solvable finite group. If \( G = 1 \), the result is trivial, so we may also assume that \( |G| > 1 \). We may also assume that \( m > 1 \) and \( n > 1 \).

**Case 1.** \( G \) contains a normal subgroup \( N \triangleleft G \) of order \( m_1n_1 \) with \( m_1 | m \) and \( n_1 | n \) satisfying \( n_1 < n \):
(1): As \( N \) is normal, \( G/N \) is solvable and \( |G/N| = \frac{m}{m_1} \) with \( \frac{m}{m_1} \) relatively prime. Since \( n_1 < n \), by the Correspondence Principle, there exists a subgroup \( N \subset H \subset G \) with the group \( H/N \) a Hall \( \frac{m}{m_1} \)-subgroup of \( G/N \) by induction. As \( |H| = \frac{m}{m_1}|N| = mn_1 < mn \) and \( H \) is also solvable, \( H \) (hence \( G \)) contains a Hall \( m \)-subgroup.

(2): Suppose that \( H' \) is another Hall \( m \)-subgroup. We must show that \( H \) and \( H' \) are conjugate. The subgroup \( HN \) is of order \( \frac{|N||H|}{|N\cap H|} = \frac{|N||H|}{|N|} = mn_1 \). Similarly, the subgroup \( H'N \) of \( G \) is of order \( n_1 \). The subgroups \( HN/N \) and \( H'N/N \) of \( G/N \) are of order \( \frac{mn_1}{m_1n_1} = \frac{m}{m_1} \), so Hall \( \frac{m}{m_1} \)-subgroups of solvable \( G/N \). By induction on \( |G| \), they are conjugate. Let \( \pi : G \to G/N \) be the canonical group epimorphism. Then there exists an \( x \) in \( G \) satisfying \( \pi(HN/N)x^{-1} = H'N/N \). It follows that \( xHNx^{-1} = H'N \). Since \( H' \) and \( xHx^{-1} \) are Hall \( m \)-subgroups of solvable \( H'N \), they are conjugate by induction.

(3): Let \( H_0 \subset G \) be a subgroup of order \( m_0 \) with \( m_0 \mid m \). Then \( |H_0N/N| \mid \frac{m}{m_1} \), so \( H_0N/N \) lies in a Hall \( \frac{m}{m_1} \)-subgroup of \( H_0N/N \) by induction. Therefore, \( H_0 \) is a Hall \( m \)-subgroup of some subgroup \( K \) of \( G \) of order \( mn_1 \). By induction on \( \frac{m}{m_1} \), the group \( K \), hence \( H_0 \), lies in a Hall \( m \)-subgroup of \( G \).

(4): By the proof above, \( h_m \) is a product of the number of Hall \( \frac{m}{m_1} \)-subgroups of \( G/N \) and the number of conjugates of \( H \) in \( HN \). The maximal power of the primes factors \( p_i^{e_i} \) of \( |HN| \) divide \( p_i^{e_i} \) and the maximal power of the primes factors of \( |G/N| \) are \( p_i^{e_i} \) for some of the \( p_i, 1 = 1, \ldots, r \). By induction, \( h_m \) is a product of two factors that satisfy condition (4). Therefore condition (4) is established for \( G \). This establishes Case 1.

**Reduction.** By Case 1, if there exists a normal subgroup \( 1 < N < G \) with \( n \nmid |N| \), we are done. So we may assume that \( n \mid |N| \) for every normal subgroup \( 1 < N < G \). In particular, this applies to each minimal normal subgroup \( H \) of \( G \). By Corollary 23.6, \( H \) is an abelian group of order \( p^r \) for some prime \( p \) and \( r \in \mathbb{Z}^+ \). Fix such an \( H \). As \( n \mid |H| \), we must have \( n = p^r \). Therefore, \( H \) is a normal Sylow \( p \)-subgroup of \( G \), hence the unique Sylow \( p \)-subgroup of \( G \). In particular, if \( N \) is any proper normal subgroup of \( G \), then \( H \subset N \). It follows that \( H \) is the unique minimal normal subgroup of \( G \).

**Case 2.** Case 1 does not hold. In particular, the group \( H \) in the Reduction is the unique minimal normal subgroup of \( G \):

(1): The non-trivial finite group \( G/H \) contains a minimal normal subgroup. By the Correspondence Principle, it follows that there exists a normal subgroup \( 1 < K < G \) satisfying \( K/H < G/H \) is a minimal
normal subgroup. By Corollary 23.6, $|K/H| = q^s$ for some prime $q$ and $s \in \mathbb{Z}^+$ with $q \neq p$. Therefore, $|K| = p^r q^s$. Let $Q \in \text{Syl}_q(G)$. Then we have $K = HQ$, since they are groups of the same order, and $Q \cap H = 1$, since $|Q|$ and $|H|$ are relatively prime.

**Claim.** $N_G(Q)$ is a Hall $m$-subgroup of $G$:

We have $N_K(Q) = K \cap N_G(Q)$. As $H \triangleleft K$ and $K \triangleleft G$, by the Frattini Argument (Exercise 22.15(??)), we know that 

$$K = HN_K(Q) \quad \text{and} \quad G = KN_G(Q).$$

As 

$$G/H = KN_G(Q)/K \cong N_G(Q)/(K \cap N_G(Q)) = N_G(Q)/N_K(Q),$$

we have $|N_G(Q)| = |G||N_K(Q)|/|K|$. As $|K| = |HN_K(Q)| = |H| |N_K(Q)|/[H:N_K(Q)]$, we have 

$$|N_G(Q)| = \frac{|G|}{|K|}|N_K(Q)| = \frac{|G|}{|H|}|H \cap N_K(Q)| = m|H \cap N_K(Q)|,$$

since $n = p^r$. To show the claim, we must show that $|H \cap N_K(Q)| = 1$. To do this we first show that $H \cap N_K(Q) \subset Z(K)$; and then we shall show that $Z(K) = 1$. Let $z \in H \cap N_K(Q)$. As $K = HQ$, there exist $h \in H$, $y \in Q$ such that $z = hy$. As $H$ is abelian, $zh = hz$. Since $z \in N_K(Q)$ and $H$ is normal in $G$, we have $(z^{-1}y^{-1}z)y = z^{-1}(y^{-1}zy) \in Q \cap H = 1$. Therefore, $z \in Z(K)$. Now we show that $Z(K) = 1$. Since $Z(K)$ is a characteristic subgroup of $K$ and $K$ is normal in $G$, $Z(K)$ is normal in $G$ by Exercise 11.9(20). If $1 < Z(K)$, then $Z(K)$ contains a minimal normal subgroup of $G$. But then $H \subset Z(K)$ by assumption. As $K = HQ$ and $H \subset Z(K)$, it follows that $Q \triangleleft K$. This then implies that $H \subset Q$. As $q \neq p$, this is a contradiction. Therefore, we have $H \cap N_K(Q) = 1$ and the claim is proven.

(2) and (4): We must show that if $A$ is another Hall $m$-subgroup of $G$, then $A$ and $N_G(Q)$ are conjugate. Let $K$ be as above, then $m \mid |AK|$ and $p^r q^s = |K| \mid |AK|$, so we have $|AK| = |G| = mn$. It follows that the subgroup $AK = G$. Therefore, $G/K = AK/K \cong A/(A \cap K)$ and $|A \cap K| = q^s$. So $A \cap K$ is a Sylow $q$-subgroup of $K$. Hence $A \cap K$ and $Q$ are conjugate in $K$ by the Second Sylow Theorem. By Lemma 23.7, $N_G(Q)$ and $N_G(A \cap K)$ are conjugate in $G$. By the claim, 

$$m = |N_G(Q)| = |N_G(A \cap K)|.$$ Since $A \cap K \triangleleft A$, we have $A \subset N_G(A \cap K)$. It follows that $A = N_G(A \cap K)$, as both are of order $m$. Consequently, $N_G(Q)$ and $A$ are conjugate subgroups of $G$. This proves (2). Since the $p^r$ conjugate subgroups of order $m$ constitute all subgroups of order $m$, (4) follows.
(3): Let \( A_0 \) be a subgroup of \( G \) of order \( m' \) with \( m' \mid m \) and \( A \) be a Hall \( m \)-subgroup of \( G \). We have \( A \cap A_0 H \) is a subset of \( A_0 H \) and \( G = AH \). By Dedekind’s Modular Law 17.7, we have

\[
A_0 H = A_0 H \cap AH = (A_0 H \cap A) H,
\]

as \( H \subset A_0 H \). Therefore, \( |A_0 H \cap AH| = |(A_0 H \cap A) H| \). It follows that \( |A_0 H \cap H| = |A_0| = m' \). By property (2) for \( A \), we have \( A_0 H \cap H \) is conjugate to \( A_0 \), so a Hall \( m \)-subgroup. \( \square \)

Remark 23.9. The conclusions of the theorem for a finite group are usually false. For example, \( A_5 \) is simple [by Abel’s Theorem – to be proven in Theorem 24.13]. It can be shown that \( A_5 \) has no subgroup of order 15 violating (1) and contains a subgroup of order 6 not contained in a subgroup of order 12, violating (3). Also \( |\text{Syl}_5(A_5)| = 6 = 3 \cdot 2 \) which violates (4). Condition (2) is also violated in general. Let \( H = \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2\mathbb{Z} \). It can be shown that \( G = \text{Aut}(H) \) is a simple group of order 168 permuting the seven groups of order two in \( H \) transitively and also the seven subgroups of order four in \( H \) transitively. Therefore, \( G \) has two distinct conjugate sets of subgroups of index 7 and order 24, so violates property (2).

An interesting special case of this theorem needs the following definition.

Definition 23.10. Let \( G \) be a finite group of order \( p^n \) with \( p \) a prime not dividing \( n \). A Hall \( n \)-subgroup of \( G \), if it exists, is called a \( p \)-complement of \( G \). We also call such a \( p \)-complement (if it exists) a \( p' \)-subgroup of \( G \). So, by definition, a \( p' \)-group if it exists, is a Hall group of order \( \frac{|G|}{p^n} \).

Corollary 23.11. If \( G \) is a finite solvable group, \( p \) a prime dividing \( |G| \), then \( G \) has a \( p \)-complement.

We now turn to another characterization of solvable groups. This will follow by proving the converse to Corollary 23.11. This is also due to Philip Hall and is the second key result in this section. To prove the theorem, we need Burnside’s Theorem 111.6 that any group of order \( p^aq^b \), \( p, q \) primes, is solvable. The proof of this theorem needs the representation theory of groups that we shall do later. We shall, however, assume its validity here.

Theorem 23.12. Let \( G \) be a non-trivial finite group. Suppose for every prime \( p \) dividing \( |G| \) that \( G \) has a \( p \)-complement. Then \( G \) is solvable.
PROOF. Suppose the result is false. We may assume that \( G \) has been chosen so that \( G \) is a counterexample of minimal order. We first show that such a \( G \) must be a simple group. Suppose that \( G \) is not a simple group, then there exists a normal subgroup \( N \) of \( G \) with \( 1 < N < G \). Let \( p \) be a prime dividing \( |G| \) and \( H \) a \( p' \)-subgroup. Then \( H \cap N \) and \( HN/N \) are \( p' \)-subgroups of \( N \) and \( G/N \), respectively. By the minimality of \( |G| \), the groups \( N \) and \( G/N \) are solvable. Hence \( G \) is solvable by Theorem 17.3, a contradiction. Therefore, we can assume that \( G \) is simple. We show that this leads to a contradiction.

Let \( |G| = p_1^{e_1} \cdots p_r^{e_r} \) with \( p_1, \ldots, p_r \) distinct primes and \( e_i > 0 \) for \( i = 1, \ldots, r \). By Burnside’s Theorem 111.6, we have \( r > 2 \). By assumption, there exist \( p_i \)-subgroups \( G_i \) of \( G \) for \( i = 1, \ldots, r \). Let \( H = G_3 \cap \cdots \cap G_r \). Then we have \( [G : G_i] = p_i^{e_i} \), hence \( [G : H] = p_3^{e_3} \cdots p_r^{e_r} \) and \( H = p_1^{e_1} p_2^{e_2} \) by Exercise 5(5). In particular, \( H \) is solvable by Burnside’s Theorem 111.6. Let \( M \) be a minimal normal subgroup of \( H \), so either a \( p_1 \)- or \( p_2 \)-subgroup, say a \( p_1 \)-subgroup. Since \( [G : H \cap G_2] = p_2^{e_2} \) \( \), we have \( [H \cap G_2] = p_1^{e_1} \), i.e., \( H \cap G_2 \) is a Sylow \( p_1 \)-subgroup of \( H \). It follows that \( M \subset H \cap G_2 \subset G_2 \). The same counting argument shows that \( |H \cap G_1| = p_2^{e_2} \). Therefore, we have \( G = (H \cap G_1)G_2 \) as sets. Since \( M \) is normal in \( H \), we have

\[
N_G(M) = \{ g^{-1}xg \mid x \in M, g \in G \} = \{(hg_2)^{-1}xh g_2 \mid x \in M, h \in H \cap G_1, g_2 \in G_2 \} = \{(g_2)^{-1}yg_2 \mid y \in M, g_2 \in G_2 \} \subset G_2.
\]

Therefore, \( N_G(M) \) is a proper non-trivial normal subgroup of \( G \). This contradicts \( G \) being simple. \( \square \)

Putting all this together, we have:

**Theorem 23.13.** Let \( G \) be a finite solvable group. Then the following are equivalent:

1. \( G \) is solvable.
2. Let \( n \) be a positive integer satisfying \( n \mid |G| \) with \( n \) relatively prime to \( |G|/n \). Then \( G \) contains a Hall \( n \)-subgroup.
3. \( G \) has a \( p \)-complement for every prime \( p \) dividing the order of \( G \).

Recall from Exercise 13.7(??), that a group \( G \) is called an (internal) semidirect product of a normal subgroup \( N \) of \( G \) and a subgroup \( H \) of \( G \) if \( G = NH \) and \( N \cap H = 1 \). [This is isomorphic to the (external) semidirect product of Exercise 11.9(14).] This means that the group homomorphism \( \theta : H \rightarrow \text{Aut}(N) \) by \( h \mapsto \theta_h \) (conjugation by \( h \) on \( N \)), satisfies \( nhn'h' = n\theta_h(n')hh' \).
Therefore, we have

**Corollary 23.14.** Let $G$ be a finite solvable group. Suppose that there exists a normal Hall $n$-subgroup of $G$ with $n$ coprime to $\frac{|G|}{n}$. Then $G$ is the semidirect product of $N$ and a Hall $\frac{|G|}{n}$-subgroup of $G$.

One can ask whether there are other interesting results concerning the existence of Hall subgroups.

**Remark 23.15.** Let $G$ be a finite group of order $mn$ with $m$ and $n$ relatively prime. The Schur-Zassenhaus Theorem is a fundamental result in group theory. It says if there exists a normal Hall $m$-subgroup $K$ of $G$, then there exists a Hall $n$-subgroup $H$ of $G$. In particular, $G$ is a semidirect product of $K$ and $H$. The Schur-Zassenhaus Theorem uses group cohomology in its proof, so will not be done here. Using the Feit-Thompson Theorem that finite groups of odd order are solvable, one then also proves that any two Hall $n$-subgroups in the Schur-Zassenhaus Theorem are conjugate and any subgroup of $H$ of $G$ of order dividing $n$ is a subgroup of some Hall $n$-subgroup of $G$.

**Exercises 23.16.**

1. Prove Lemma 23.7
2. (Wielandt) Suppose that $G$ is a finite group having three solvable groups $H_1$, $H_2$, and $H_3$ satisfying $[G : H_1]$, $[G : H_2]$, $[G : H_3]$ are pairwise relatively prime. Prove that $G$ is solvable.
3. Let $G$ be a finite group. Suppose every non-trivial subgroup $H$ of $G$ has a subgroup of index $p$ for every prime dividing $G$. Show that $G$ is solvable.
4. Let $G$ be a group of order 12. Suppose that $G$ is not isomorphic to $A_4$. Show that $G$ contains a normal Sylow 2-subgroup or a normal Sylow 3-subgroup. Classify all groups of order 12 up to isomorphism.

24. The Symmetric and Alternating Groups

Recall if $\sigma \in S_n$, we write $\sigma$ as

$$
\begin{pmatrix}
1 & 2 & \ldots & n \\
\sigma(1) & \sigma(2) & \ldots & \sigma(n)
\end{pmatrix}
$$

where the top row is the elements of the domain and the bottom row the corresponding values. For example, if

$$
\sigma = \begin{pmatrix} 1 & 2 & 3 & 4 \\ 2 & 1 & 4 & 3 \end{pmatrix} \quad \text{and} \quad \tau = \begin{pmatrix} 1 & 2 & 3 & 4 \\ 1 & 3 & 2 & 4 \end{pmatrix}
$$

then

$$
\tau \sigma = \begin{pmatrix} 1 & 2 & 3 & 4 \\ 3 & 1 & 4 & 2 \end{pmatrix} \quad \text{and} \quad \sigma \tau = \begin{pmatrix} 1 & 2 & 3 & 4 \\ 2 & 4 & 1 & 3 \end{pmatrix}.
$$
We want to write permutations in other ways.

**Definition 24.1.** A permutation \( \alpha \in S_n \) is called a *cycle of length* \( r \), or simply an *r-cycle*, if there exist distinct elements \( a_1, \ldots, a_r \) in \( \{1, \ldots, n\} \) (so \( r \leq n \)) satisfying

1. \( \alpha(a_i) = a_{i+1} \) for \( i = 1, \ldots, r - 1 \).
2. \( \alpha(a_r) = a_1 \).
3. \( \alpha(a) = a \) for all \( a \in \{1, \ldots, n\} \setminus \{a_1, \ldots, a_r\} \).

We denote the above \( r \)-cycle by \((a_1 \cdots a_r)\). If \( r > 1 \), we say that \( \alpha \) is *non-trivial*.

**Examples 24.2.**

(i) \((1)\) \((a_1) = 1\) the identity on \( \{1, \ldots, n\} \).

(ii) \[
\begin{pmatrix}
1 & 2 & 3 & 4 \\
1 & 3 & 2 & 4
\end{pmatrix}
= (23).
\]

(iii) \[
\begin{pmatrix}
1 & 2 & 3 & 4 \\
2 & 1 & 4 & 3
\end{pmatrix}
= (34)(12) = (12)(34).
\]

(iv) \[
\begin{pmatrix}
1 & 2 & 3 & 4 \\
3 & 1 & 4 & 2
\end{pmatrix}
= (1342).
\]

(v) \[
\begin{pmatrix}
1 & 2 & 3 & 4 \\
2 & 4 & 1 & 3
\end{pmatrix}
= (1243).
\]

The basic properties of cycles are given by:

**Properties 24.3.** Let \( \alpha = (a_1 \cdots a_r) \) in \( S_n \) with \( n \geq 2 \). Then

1. \( \alpha = (a_i a_{i+1} \cdots a_r a_1 \cdots a_{i-1}) \) if \( 1 \leq i \leq r \).
2. \( \alpha \) has order \( r \).
3. \( \alpha^{-1} = (a_r \cdots a_1) \).
4. If \( \sigma \in S_n \), then \( \sigma \alpha \sigma^{-1} = (\sigma(a_1) \cdots \sigma(a_r)) \).

**Proof.** We leave the proof of (1) and (3) as an exercise.

(2): \( \alpha^i(a_1) = a_{i+1} \neq a_1 \) for \( 1 \leq i < r \), so the order of \( \alpha \) is at least \( r \). It is easy to check that \( \alpha^i(a_i) = a_i \) for \( 1 \leq i \leq r \), so the order of \( \alpha = r \).

(4): we have

\[
\sigma \alpha \sigma^{-1}(j) = \begin{cases} 
\sigma(a_{i+1}) & \text{if } \sigma^{-1}(j) = a_i, \text{ i.e., } j = \sigma(a_i) \text{ for } i < r. \\
\sigma(a_1) & \text{if } \sigma^{-1}(j) = a_r, \text{ i.e., } j = \sigma(a_r). \\
j (= \sigma \sigma^{-1}(j)) & \text{otherwise.}
\end{cases}
\]

We say that two cycles \( \alpha = (a_1 \cdots a_r) \) and \( \beta = (b_1 \cdots b_l) \) are *disjoint* if \( \{a_1, \ldots, a_r\} \cap \{b_1, \ldots, b_l\} = \emptyset \).
Lemma 24.4. If $\alpha$ and $\beta$ are disjoint cycles, then $\alpha$ and $\beta$ commute, i.e., $\alpha \beta = \beta \alpha$.

**Proof.** $\beta(a_i) = a_i$ for all $i$ and $\alpha(b_j) = b_j$ for all $j$. \qed

We now use the evaluation action in Example 21.20, which we recall.

**Construction 24.5.** Suppose that $S = \{1, \ldots, n\}$ and $\gamma \in S_n$ with $n > 1$. Let $\Gamma = \langle \gamma \rangle \subset S_n$. Then $S$ is a $\Gamma$-set via the evaluation action $* : \Gamma \times S \to S$ defined by $\gamma^j * a = \gamma^j(a)$.

The orbit of $a \in S$ is $\Gamma \ast a = \{\gamma^j(a) \mid j \in \mathbb{Z}\}$. Since $|\Gamma| \leq |S_n| = n! < \infty$, just as in the proof of the Classification of Cyclic Groups, there exists a least positive integer $m = m(a)$ satisfying $\gamma^m(a) = a$. Therefore, the orbit of $a$ is $\Gamma \ast a = \{a, \gamma(a), \ldots, \gamma^{m-1}(a)\} \subset S$. Associate to this orbit, the $m$-cycle $\gamma_a = (\gamma(a) \cdots \gamma^{m-1}(a))$ in $S_n$.

Let $O$ be a system of representatives for the equivalence relation $\sim_{\Gamma}$ arising from the $\Gamma$-action on $S$ and let $b \in \{1, \ldots, n\}$. There exists an $a \in O$ satisfying $b \in \Gamma \ast a$ and $\gamma(b) = \gamma_a(b)$. By the Mantra of $G_p$-actions, $S = \bigvee_O \Gamma \ast a$. Since disjoint cycles commute, we have $\gamma = \prod_O \gamma_a$, i.e., $\gamma$ is a product of disjoint cycles with the product unique up to order. We call this the **(full) cycle decomposition** of $\gamma$.

For example,

$$\begin{pmatrix} 1 & 2 & 3 & 4 & 5 & 6 \\ 3 & 2 & 5 & 4 & 1 & 6 \end{pmatrix} = (135)(2)(4)(6)$$

and

$$\begin{pmatrix} 1 & 2 & 3 & 4 & 5 & 6 \\ 2 & 1 & 5 & 3 & 4 & 6 \end{pmatrix} = (12)(354)(6).$$

Note (again) that a 1-cycle $(a)$ is the identity 1s. Therefore, the **fixed points** of the $\Gamma$-action (i.e., elements $a$ in $S$ satisfying $\gamma(a) = a$) are precisely the elements of $S$ occurring in the 1-cycles (if any) in a cycle decomposition of $\gamma$. We usually delete 1-cycles in the cycle decomposition of a permutation, e.g.,

$$\begin{pmatrix} 1 & 2 & 3 & 4 & 5 & 6 \\ 3 & 2 & 5 & 4 & 1 & 6 \end{pmatrix}$$

is written (135). We shall call the cycle decomposition **full** if we do not delete the 1-cycles.

By Properties 24.3(4), we have the following useful observation:

**Remark 24.6.** If $\prod_O \gamma_a$ is a cycle decomposition for $\gamma$ in $S_n$, then $\prod_O \sigma \gamma_a \sigma^{-1}$ is a cycle decomposition for $\sigma \gamma \sigma^{-1}$. (Cf. Exercise 24.24(4).)

We turn to another representation of permutations. We call a 2-cycle a **transposition**. Note if $i \neq j$ then $(ij)^{-1} = (ij)$ has order 2 as needed. We first show that transpositions generate $S_n$ (for $n > 1$).
Proposition 24.7. Let \( n > 1 \), then every element in \( S_n \) is a product of the transpositions \((12), (13), \ldots, (1n)\).

Proof. We first show that \( S_n \) is generated by transpositions. If \( \sigma = (a_1 \ldots a_r) \), then
\[
\sigma = (a_r^{-1}a_r) \cdots (a_2a_r)(a_1a_r)
\]
is a product of \( r - 1 \) transpositions.

To show that \( S_n \) is generated by the specific transpositions in the proposition, we use Properties 24.3(4) and (5), respectively, to see that
\[
(1i) = (i1) = (i1)^{-1} \quad \text{if } i \neq 1
\]
\[
(ij) = (1i)(1j)(1i)^{-1} = (1i)(1j)(1i) \quad \text{if } 1, i, \text{ and } j \text{ are distinct.}
\]
The result follows using the cycle decomposition of permutations. \( \square \)

Remarks 24.9. If \( \sigma \in S_n \), with \( n > 1 \), then \( \sigma \) is not a product of a unique number of transpositions. However, by Equation 11.6, we have group homomorphisms
\[
S_n \xrightarrow{\theta} \text{Perm}_n(\mathbb{R}) \xrightarrow{\det} \{\pm 1\},
\]
given by \( \sigma \mapsto [T_{\sigma}]_{S_n} \mapsto \det[T_{\sigma}]_{S_n} \), where \( S_n \) is the standard basis for \( \mathbb{R}^n \). The first map is an isomorphism, the second an epimorphism. So we have the alternating group \( A_n := \ker \det \circ \theta \circ S_n \) and \( S_n/A_n \cong \{\pm 1\} \cong \mathbb{Z}/2\mathbb{Z} \) by the First Isomorphism Theorem. In particular, we have \( |S_n : A_n| = 2 \). If \( \sigma \) is a transposition, \( \det[T_{\sigma}]_{S_n} = -1 \) (as it corresponds to interchanging two columns of the identity matrix). It follows that if \( \tau \in S_n \) is a product of \( r \) transpositions and a product of \( s \) transpositions, then \((-1)^r = (-1)^s \), i.e., \( r \equiv s \mod 2 \). It follows by Equation 24.8, if \( \sigma \) is an \( r \)-cycle with \( r \) odd, then \( \sigma \in A_n \) and if \( r \) is even, then \( \sigma \in S_n \setminus A_n \). In particular, if \( n \geq 3 \), then every 3-cycle lies in \( A_n \). Using the cycle decomposition of any permutation, we conclude:
\[
A_n = \{ \sigma \in S_n \mid \sigma \text{ is a product of an even number of transpositions} \}
\]
and
\[
S_n \setminus A_n = \{ \sigma \in S_n \mid \sigma \text{ is a product of an odd number of transpositions} \}.
\]
Elements in \( A_n \) are called even permutation and elements in \( S_n \setminus A_n \) are called odd permutations. So if \( \tau \) is an odd permutation, we have \( S_n = A_n \vee \tau A_n \).

Actually, one must be a bit careful in the above so that we do not get a circular argument. This would depend on the proof of the determinant, whose existence we have not proven. One can prove the properties of the determinant so that one does not have a circular
argument, the problem being that interchanging two rows of a square matrix changes the sign of its determinant. [Cf. Section 118 for a sophisticated proof of the determinant and its properties.]

The possible place for the circular reasoning arises from the definition of the signum map. For completeness, we shall give a proof for this. The standard proof uses polynomials, and although simple seems out of place, so we give a group theoretic proof. Let $\sigma$ be a permutation in $S_n$ and suppose that $\sigma = \gamma_1 \cdots \gamma_r$ is its full cycle decomposition. Define the signum of $\sigma$ by $\text{sgn}(\sigma) = (-1)^{n-r}$. As a cycle decomposition is unique, this defines a map $\text{sgn} : S_n \rightarrow \{\pm 1\}$.

Note if $\sigma$ is an $s$-cycle, then in its full cycle decomposition, the number of 1-cycles is $n-s$, so $\text{sgn} \sigma = (-1)^{s-1}$. In particular, if $\sigma$ is a transposition, then $\text{sgn}(\sigma) = -1$. The result that we need is that $\text{sgn}$ is a group homomorphism. Clearly, the $\text{sgn}$ of disjoint cycles is the product of the $\text{sgn}$ of those cycles. As every permutation is a product of transpositions, the result now follows easily from the following:

Let $a, b, c_1, \ldots, c_r, d_1, \ldots, d_s$ be distinct elements in $\{1, \ldots, n\}$ and $\tau = (ab)$. If $\sigma = (ac_1 \cdots c_r)$, then $\tau \sigma = (c_1 \cdots c_r ba)$, hence

$$\text{sgn}(\tau \sigma) = (-1)^{r+1} = -\text{sgn}(\sigma);$$

and if $\sigma = (ac_1 \cdots c_r bd_1 \cdots d_s)$, then $\tau \sigma = (ac_1 \cdots c_r)(bd_1 \cdots d_s)$, hence

$$\text{sgn}(\tau \sigma) = (-1)^{r+s} = -(-1)^{r+s+1} = -\text{sgn}(\sigma).$$

In particular, the parity of the number of transpositions in a decomposition of a permutation is invariant.

We turn to the study the alternating group. We know that 3-cycles are even. We start by showing that 3-cycles generate $A_n$. As $A_n$, $n = 1, 2$ is the trivial group, it is generated by the empty set of 3-cycles, so we may assume that $n \geq 3$.

**Proposition 24.10.** The alternating group $A_n$ is a normal subgroup of index two in $S_n$, and, if $n \geq 3$, then $A_n$ is generated by the 3-cycles:

$$(123), (124), \ldots, (12n).$$

**Proof.** We need only show the last statement. We first show that $A_n$ is generated by 3-cycles. We have shown that every element in $A_n$ is a product of an even number of transpositions, so it suffices to show a product of two distinct transpositions is a product of 3-cycles. Suppose that $i, j, k, l$ are distinct elements in $\{1, \ldots, n\}$. Then this follows by the equations

$$(ij)(ik) = (ikj)$$

$$(ij)(kl) = (kil)(ijk).$$
To finish, we show that the alleged 3-cycles generate. Let $i, j, k$ be distinct elements in $\{1, \ldots, n\}$, and different from 1 and 2. Then, using Property 24.3, this follows from the equations:

\[
(ijk) = (12i)(2jk)(12i)^{-1}
\]
\[
(2jk) = (12j)(12k)(12j)^{-1}
\]
\[
(1jk) = (12k)^{-1}(12j)(12k).
\]

□

Remark 24.11. We make the following remark to simplify notation below. Let $n \geq 3$ and $K \subset A_n$ a subgroup. If $\sigma \in S_n$, then $\sigma K \sigma^{-1} \subset \sigma A_n \sigma = A_n \triangleleft S_n$. In particular, $K$ contains a 3-cycle if and only if $\sigma K \sigma^{-1}$ does.

The key to Abel’s Theorem on the simplicity of the alternating group $A_n$ for $n \geq 5$, is the following:

Lemma 24.12. Let $K$ be a normal subgroup of $A_n$. If $K$ contains a 3-cycle, then $K = A_n$.

Proof. As $K$ contains a 3-cycle, we have $n \geq 3$. By the proposition, it suffices to show that $(12k)$ lies in $K$ for $k = 3, \ldots, n$. By Remark 24.11, we may assume that $(123)$ lies in $K$, hence $(213) = (123)^{-1}$ lies in $K$. Let $\sigma = (12)(3k)$ with $k > 3$, a product of an even number of transpositions, so an element of $A_n$. As $K \triangleleft A_n$, we have

\[
(12k) = \sigma(213)\sigma^{-1}
\]

lies in $K$, for all $k \geq 3$

as needed. □

We now prove Abel’s Theorem. This gives our first examples of non-abelian simple groups.

Theorem 24.13. (Abel’s Theorem) Let $n \neq 4$, then the alternating group $A_n$ is simple.

Proof. We know that $A_2$ and $A_3$ are simple, so assume that $n \geq 5$, and $1 < K \triangleleft A_n$. We must show $K = A_n$. By the lemma, it suffices to show that $K$ contains a (single) 3-cycle.

Let $e \neq \alpha \in K$ be chosen such that

\[
\alpha \text{ moves precisely } m \text{ elements in } \{1, \ldots, n\}
\]

and

\[
\text{no other non-identity element in } K \text{ moves a fewer number of elements},
\]
where we say $\alpha$ moves $i$, if $\alpha(i) \neq i$, i.e.,

$$m$$ is the minimal number of elements moved by any non-identity element in $K$.

Let

$$\alpha = \gamma_1 \cdots \gamma_r$$

be the decomposition of $\alpha$ into a product of disjoint cycles. We may assume that $\gamma_1$ is a $k$-cycle with $k$ maximal among the lengths of $\gamma_1, \ldots, \gamma_r$.

As $\alpha$ is not the identity, we must have $m \geq 2$. If $m = 2$, then we must have $\alpha = \gamma_1$ and be a transposition, hence it is an odd permutation, so not an element of $A_n$, a contradiction. If $m = 3$, then we must have $\alpha = \gamma_1$ and be a 3-cycle, and we are done by the lemma. So we may assume that $m \geq 4$. For clarity, it is convenient to changing notation if necessary. So, using Remark 24.11, we may assume that

$$\alpha$$ moves 1, 2, 3, 4.

$$\alpha$$ moves 5 if $m \geq 5$.

$$\gamma_1 = (12 \cdots k).$$

We have two cases:

**Case 1.** $k \geq 3$: If this is the case, then $m \geq 5$:

For if not, then we have $m = 4$, and we must have $\alpha = \gamma_1 = (1234)$ and be a 4-cycle which is not in $A_n$, a contradiction.

**Case 2.** $k = 2$: Then each $\gamma_i$ is a transposition or a 1-cycle, and we may assume, using Remark 24.11, that $\alpha = (12)(34) \cdots$:

Set $\beta = (345)$ in $A_n$. As $K \triangleleft A_n$, we have $\alpha_1 = \beta \alpha \beta^{-1}$ lies in $K$ and

$$\alpha_1 = \beta \alpha \beta^{-1} = \begin{cases} (124 \cdots) \cdots & \text{in Case 1.} \\ (12)(45) \cdots & \text{in Case 2.} \end{cases}$$

In either case, we check that $\alpha(3) \neq \alpha_1(3)$ (e.g., in Case 1, if $k = 3$, then $\alpha_1(3) \neq 1$ and if $k \geq 4$, then $\alpha_1(3) \neq 4$), so $e \neq \alpha^{-1} \alpha_1$ lies in $K$; and $\alpha(1) = \alpha_1(1)$, so $\alpha^{-1} \alpha_1$ fixes 1. Since $\beta$ fixes every $i > 5$, we have $\alpha^{-1} \alpha_1$ fixes every $i > 5$ that $\alpha$ fixes. As $m \geq 5$ in Case 1 and in this case $\alpha^{-1} \alpha_1$ fixes 1, it follows by the minimality of $m$ that this can happen only if $\alpha$ satisfies Case 2. But then $\alpha(1) = \alpha_1(1)$ and $\alpha(2) = \alpha_1(2)$, so $\alpha^{-1} \alpha_1$ fixes 1 and 2. As $m \geq 4$ (with $\alpha$ fixing 5 if $m = 4$, but $\alpha_1$ not fixing 5), it follows that $\alpha^{-1} \alpha_1$ moves at most $m - 1 = m - 2 + 1$ elements contradicting the minimality of $m$. The result follows. \qed
Remark 24.14. If $G$ is a group containing a non-abelian simple group, then $G$ is not solvable by Theorem 17.3. Therefore, $S_n$ is not solvable for any $n \geq 5$.

We wish to show that $A_5$ is the only simple group of order 60 up to isomorphism (non-abelian as its order is not a prime). It is the first non-abelian simple group by computation, using the Sylow Theorems. We make the general observation:

Lemma 24.15. Let $n \geq 5$, then $A_n$ is the only subgroup of $S_n$ of index two.

Proof. Suppose that $K$ is a normal subgroup of $G$ of index two and $K \neq A_n$. By the Second Isomorphism Theorem, we conclude $A_nK$ is a group with $K < A_nK$ normal and $A_nK/K \cong A_n/(A_n \cap K)$. As $A_n < A_nK \subseteq S_n$, we must have $S_n = A_nK$, so $2 = [A_n : A_n \cap K]$. In particular, $1 < A_n \cap K < A_n$. It follows that $A_n \cap K$ is a proper normal subgroup of the simple group $A_n$, a contradiction. □

Theorem 24.16. Up to isomorphism, the alternating group $A_5$ is the only simple group of order 60.

Proof. Let $G$ be a simple group of order 60. We prove this in three steps.

Step 1. If $G$ contains a subgroup $H$ of index 5, then $G \cong A_5$:
Let $\lambda : G \to \Sigma(G/H)$ given by $x \mapsto (\lambda_x : gH \mapsto xgH)$ be the homomorphism defined in the General Cayley Theorem 12.4. Since $H < G$, we have $\ker(\lambda) \subset H < G$. Consequently, $\lambda$ must be monic (since not the trivial map), as $G$ is simple. Therefore, $G$ is isomorphic to $\lambda(G) \subset \Sigma(G/H)$, a subgroup of index two. By the lemma $G \cong A_5$.

Step 2. $G$ does not contain a subgroup of index three:
Suppose that $H$ is a subgroup of index three in $G$. Then $|G| \nmid |G : H|!$, which implies that $G$ is not simple by Useful Counting 12.7.

Step 3. Finish:
By the first two steps, we may assume that $G$ has no subgroup of index 3 or index 5. Let $P$ be a Sylow 2-subgroup of $G$. Then $P$ is not a normal subgroup of $G$ by assumption, hence by the Third Sylow Theorem, we have $[G : N_G(P)] = 1 + 2k \mid 3 \cdot 5$, for some integer $k \geq 1$. By the first two steps, we must have $|G : N_G(P)| = 15$, so $|N_G(P)| = 4$. (Note that this means that $P = N_G(P)$.) Let $Q \neq P$ be another Sylow 2-subgroup of $G$, and set $I = P \cap Q$. We may also assume that $Q$ and $P$ were chosen such that $I = P \cap Q$ has maximal cardinality.

Case 1. $1 < I$, so $|I| = 2$:
In this case, we have, by Exercise 10.16(4) (cf. Equation 13.5),
\[ |PQ| = |P| |Q|/|P \cap Q| = |P| |Q|/|I| = 2^3. \]
As both \(P\) and \(Q\) are groups of order \(2^2\), they are both abelian, so normalize \(I\). Consequently, \(P, Q \subseteq N_G(I)\). It follows that
\[ 2^3 = |PQ| \leq |N_G(I)| \cdot 2^2 \cdot 3 \cdot 5 \quad \text{and} \quad 2^2 \mid |N_G(I)|. \]
Therefore, \(|N_G(I)| = 2^2 \cdot 3, 2^2 \cdot 5, \text{ or } 2^2 \cdot 3 \cdot 5.\) By the first two steps, we conclude that \(|N_G(I)| = 2^2 \cdot 3 \cdot 5 = |G|\). Hence \(I \triangleleft N_G(I) = G\), contradicting the simplicity of \(G\).

**Case 2.** \(I = 1:\)**

As \(Q\) and \(P\) were chosen with \(I = P \cap Q\) of maximal cardinality, and each Sylow 2-subgroup contains \((4 - 1) = 3\) non-identity elements, the Sylow 2-subgroups of \(G\) contain \(45 = 15(4 - 1)\) distinct non-identity elements. By the simplicity of \(G\), the number of the Sylow 5-subgroups, which divides \(2^2 \cdot 3\), must be at least six. Each Sylow 5-subgroup of \(G\) is cyclic of order 5, so \(G\) must contain at least \(6 \cdot (5 - 1) = 24\) distinct elements of order 5, which implies that \(|G| \geq 45 + 24 > 60\), a contradiction.

**Remarks 24.17.** Let \(F\) be a field, then the group
\[
\text{PSL}_n(F) := \frac{\text{SL}_n(F)}{Z(\text{SL}_n(F))}
\]
is called a *projective special linear group* of \(F\). It can be shown that \(\text{PSL}_n(F)\) is simple for all \(n \geq 3\).

We shall see later that for all primes \(p\) and \(q = p^n\), with \(n \in \mathbb{Z}^+\), there exists a field \(\mathbb{F}_q\) with \(q\) elements, and unique up to isomorphism of fields [definition?]. Computation shows that

\[
|\text{PSL}_2(\mathbb{F}_q)| = \begin{cases} 
(q + 1)(q^2 - q) & \text{if } q = 2^n \\
\frac{1}{2}(q + 1)(q^2 - q) & \text{if } q = p^n \text{ with } p \text{ and odd prime}
\end{cases}
\]
(cf. Exercise 9.12(1)). It can be shown that \(\text{PSL}_2(\mathbb{F}_q)\) is simple if and only if \(q > 3\). The computation shows that \(\text{PSL}_2(\mathbb{F}_4)\) and \(\text{PSL}_2(\mathbb{F}_5)\) both have order 60, so are isomorphic to \(A_5\). The simple group \(\text{PSL}_2(\mathbb{F}_7)\) of order 168 is the next non-abelian simple group up to isomorphism. Computation shows that the groups \(A_8\) and \(\text{PSL}_3(\mathbb{F}_4)\) both have \(8!/2\) elements. It can also be shown that they are not isomorphic, i.e., there exist non-isomorphic simple groups of the same order.

Our results and exercises have shown that groups of orders \(p^s q\) (\(s \geq 1\)), \(p^2 q^2\), and \(pqr\) are not simple when \(p, q, \text{ and } r\) are primes. It follows that these groups are therefore solvable, by induction and reduction to simpler cases. Burnside showed that groups of order \(p^aq^b\) are solvable...
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for any primes $p$ and $q$ for over a hundred years ago. The original proof
did not use abstract group theory (but group representation theory).
A group theoretic proof was finally found, but is much more difficult
than this original proof. More spectacularly, Feit-Thompson proved
that any group of odd order is solvable. The proof is very difficult,
and was the first major step in classifying finite simple groups. Groups
of even order are, therefore, much harder to understand. We give two
examples when we can say something more.

The first example, is left as an easy exercise (cf. the proof of Lemma
24.15):

**Proposition 24.18.** Let $n \geq 5$. If $H$ is a normal subgroup of $S_n$, then
either $H = 1$, $H = A_n$, or $H = S_n$.

The order of $A_5$ is divisible by 4. We show groups of even order
larger than two but not divisible by 4 are never simple, i.e., if $G$ is a
group of order $2n$, with $n > 1$ an odd integer, then $G$ is not simple.

**Review 24.19.** Let $G$ be a nontrivial finite group. We have seen that
Cayley’s Theorem gives a group monomorphism

$$
\lambda : G \to \Sigma(G) \text{ defined by } x \mapsto (\lambda_x : a \mapsto xa)
$$

called the (left) regular representation of $G$. If $xa = \lambda_x(a) = a$, then
$x = e$, i.e., the permutation $\lambda_a$ has a fixed point if and only if $a = e$.
Note also that the regular representation corresponds to the translation
action of $G$ on $S = G$ which is transitive, i.e., $S$ is the only orbit under
this action. In particular, this action has no fixed points (one-point
orbits).

We need two lemmas.

**Lemma 24.20.** Let $G$ be a group of order $mn$ and $\lambda : G \to \Sigma(G)$ the
regular representation with $m > 1$. If the element $a$ in $G$ has order $m$,
then

1. $\lambda_a$ is a product of $n$ disjoint $m$-cycles.
2. $\lambda_a$ is an odd permutation if and only if $m$ is even and $n$ is odd.

[Note. This characterizes elements in $\lambda(G)$, but not in $\Sigma(G)$.]

**Proof.** (1): By the review, we know if $1 \leq r < m$, then the
permutation $\lambda_a^r = (\lambda_a)^r$ has no fixed points. Let $\lambda_a = \gamma_1 \cdots \gamma_s$ be a
full cycle decomposition of $\lambda_a$. As $\lambda_a^r$ has no fixed points for $1 \leq r < m$,
each $\gamma_i$ must be a cycle of order at least $m$. But $\lambda_a^m = 1_G$, so each $\gamma_i$
must be a cycle of length $m$. Since each element of $G$ occurs precisely
once in the cycle decomposition for $\lambda_a$, we must have $s = n$. 
By (1), we have a full cycle decomposition of $\lambda_a$ given by $\lambda_a = \gamma_1 \cdots \gamma_n$, with each $\gamma_i$ an $m$-cycle. Therefore, $\lambda_a$ is an odd permutation if and only if $n(m - 1)$ is odd, i.e., if and only if $n$ is odd and $m$ is even. □

**Lemma 24.21.** Let $G$ be a finite group and $\lambda : G \to \Sigma(G)$ the regular representation. If $\lambda(G)$ contains an odd permutation, then there exists a normal subgroup of $G$ of index two.

**Proof.** Let $|G| = n$ and identify $\Sigma(G)$ and $S_n$. As $\lambda(G)$ contains an odd permutation, $H := A_n \cap \lambda(G) < \lambda(G)$ and $\lambda(G)A_n = \Sigma(G)$. The regular representation is monic, so, using the Second Isomorphism Theorem, we must have

$$[G : \lambda^{-1}(H)] = [\lambda(G) : H] = [\lambda(G) : A_n \cap \lambda(G)]$$

$$= [A_n \lambda(G) : A_n] = [S_n : A_n] = 2.$$

Therefore, $\lambda^{-1}(H) < G$ is a subgroup of index two hence is also normal. □

**Application 24.22.** Let $G$ be a finite group of order $2n$ with $n$ odd. Then $G$ contains a normal subgroup of index two. In particular, if $n > 1$, $G$ is not simple.

**Proof.** By Cauchy’s Theorem, there exists an element $a$ in $G$ of order two. Let $\lambda : G \to \Sigma(G)$ be the regular representation. By Lemma 24.20, the permutation $\lambda_a$ is a product of $n$ disjoint transpositions, hence is an odd permutation. The result now follows by Lemma 24.21. □

This application can be pushed a bit further. Since all Sylow $p$-subgroups of a group are conjugate, if one is abelian (respectively, cyclic) all are.

**Theorem 24.23.** Let $G$ be a finite group of order $2^r m$ with $m$ odd. If $G$ contains a cyclic Sylow 2-subgroup then there exists a normal subgroup of $G$ of index $2^r$. In particular, if $m > 1$ or $r > 1$, then $G$ is not simple.

**Proof.** We may assume that $m > 1$ and $r \geq 1$. Let $\lambda : G \to \Sigma(G)$ be the regular representation. By Lemma 24.20, the group $\lambda(G)$ contains an odd permutation, so by Lemma 24.21, there exists a normal subgroup $N$ of $G$ of index two. Hence $|N| = 2^{r-1}m$, and we may assume that $r > 1$. Let $P_0 \in \text{Syl}_2(N)$. By the Fourth Theorem, there exists $P \in \text{Syl}_2(G)$ containing $P_0$. By the remark above, $P$ is cyclic, say $P = \langle x \rangle$, so $P_0 = \langle x^2 \rangle$ by the Cyclic Subgroup Theorem. By induction on $r$, there exists a normal subgroup $H$ in $N$ of
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index $2^{r-1}$. By the Second Isomorphism Theorem, $P_0 H \subset N$ is a subgroup, and $P_0$ and $H$ have relatively prime orders, so $P_0 \cap H = 1$ (why?). It follows by counting that $N = P_0 H$. Let $H^x := xHx^{-1}$. We have $G = N \vee xN$, since $x \notin N$. Therefore, $H \triangleleft G$ if and only if $H^x = H$, as $H \triangleleft N$. If this is the case, we are done, so we may assume that $H \neq H^x$. The group $N$ being normal in $G$ implies that $N^x = xNx^{-1} = N$. Since $N = P_0 H$ and $P_0^x = xP_0x^{-1} = P_0$ in cyclic $P$, we have $N = P_0 H = N^x = P_0^x H^x$, so $H^x \subset N$. As $|H| = |H^x|$, we know that $2 \nmid |H^x|$. By the Second Isomorphism Theorem, as $H \triangleleft N$ and $H^x \subset N$, we have $HH^x$ is a subgroup of $N$. In particular, $|H| \mid |HH^x|$. But we also have $2 \nmid |H|/|H^x| = |H \cap H^x|$, so we must have $|HH^x| \leq |N|/2^e = |H|$, hence $|HH^x| = |H| = |H^x|$. We conclude that $H = H \cap H^x = H^x$, a contradiction. Therefore, $H \triangleleft G$. □

Exercises 24.24.

1. Let $\alpha$ be a product of disjoint cycles $\gamma_1, \ldots, \gamma_m$. Show that the order of $\alpha$ is the least common multiple of the orders of the $\gamma_i$, $1 \leq i \leq m$. In particular, show that this means that the order of each $\gamma_i$ divides the order of $\alpha$ and further, that if $p > 1$ is a prime, then every power of a $p$-cycle is either a $p$-cycle of the identity.

2. Let $p$ be a prime and $G$ a subgroup of the symmetric group $S_p$ containing a $p$-cycle and a transposition. Show that $G = S_p$.

3. Show that $S_4$ and $SL_2(\mathbb{Z}/3\mathbb{Z})$ are not isomorphic.

4. We say that an element $\alpha \in S_n$ has cycle type $(r_1, \ldots, r_m)$ if $\alpha$ is a product of disjoint cycles $\gamma_1, \ldots, \gamma_m$ with $\gamma_i$ an $r_i$-cycle for $1 \leq i \leq m$ and $r_1 \leq \cdots \leq r_m$. Prove that two elements of $S_n$ are conjugate in $S_n$ if and only if they have the same cycle type.

5. A permutation in $S_n$ is called regular if it is the identity or has no fixed points and is the product of disjoint cycles of the same length. Show that a permutation is regular if and only if it is a power of an $n$-cycle.

6. Let $\alpha$ be an $r$-cycle. If $k > 0$ is an integer and $(r, k)$ the gcd of $r$ and $k$, show that $\alpha^k$ is a product of $(n, k)$ disjoint cycles, each of length $r/(r, k)$.

7. Show the following:

   (a) $A_4$ can be generated by two elements $x$ and $y$ satisfying $x^2 = y^3 = (xy)^3 = 1$.

   (b) $S_4$ can be generated by two elements $x$ and $y$ satisfying $x^2 = y^3 = (xy)^4 = 1$.

   (c) $A_5$ can be generated by two elements $x$ and $y$ satisfying $x^2 = y^3 = (xy)^5 = 1$. 
8. Show that $A_5$ is generated by 5-cycles. Generalize.

9. Let $n \geq 5$. Show if $H$ is a normal subgroup of $S_n$, then either $H = 1$, $H = A_n$, or $H = S_n$.

10. Let $p$ be a prime and $H$ the cyclic subgroup generated by a $p$-cycle in $S_p$. Determine the order of the normalizer $N_{S_p}(H)$. Prove your determination is correct. Generalize this to the cyclic subgroup generated by an $n$-cycle in $S_n$.

11. Let $G$ be a finite group and $N$ a normal subgroup of $G$ of prime index $p$. Prove that $[G, G] \subset N$ and use it to prove that $S_n$, $n \geq 5$ is not polycyclic, hence not solvable without using Abel’s Theorem.

12. Let $n \geq 5$. Prove that the only proper subgroup of $S_n$ of index at most $n - 1$ is $A_n$.

13. Let $G$ be a free group on basis $\mathcal{B}$ with $\mathcal{B}$ having at least two elements. Show that $G$ is not solvable.

14. Show that the free product $\mathbb{Z}/2\mathbb{Z} \ast \mathbb{Z}/5\mathbb{Z}$ (cf. Definition 18.12 and Theorem 18.13) is not solvable.

15. Let $A_\infty := \bigcup_{i=1}^{\infty} A_n$ and $S_\infty := \bigcup_{i=1}^{\infty} S_n$ be groups defined in the obvious way (with $A_n$ and $S_n$ respective subgroups for all $n$). Show that $A_\infty$ is a normal subgroup of index two in $S_\infty$ and $1 < A_\infty < S_\infty$ is a composition series. Also show the infinite abelian subgroup of $S_\infty$ generated by the transpositions $(2k - 1, 2k)$, $k \geq 1$, does not have a composition series. (Cf. Exercise 17.14(11).)
Part 3

Ring Theory
CHAPTER V

General Properties of Rings

In this chapter, we begin our formal study of rings, especially that of commutative rings. We find the analogues in ring theory for the notions and basic theorems we learned in group theory. The analogue in ring theory of a normal subgroup is called an ideal. As in group theory, this turns out to be those subsets of a ring that are kernels of ring homomorphisms. This leads to analogous isomorphism theorems and an analogous correspondence principle.

When studying the integers, we looked at integers themselves. In this chapter, we switch our perspective from studying elements to studying ideals. We discover the type of ideal corresponding to a prime integer in commutative rings. This type of ideal, called a prime ideal, is the major object of study in commutative ring theory whose definition is based upon Euclid’s Lemma. (For elements the analogue of a prime integer is called an irreducible element.) We also introduce an important generalization of the ring of integers called a principal ideal domain or PID. Although this ring is a very special type of domain (a commutative ring in which \(ab = 0\) implies \(a = 0\) or \(b = 0\)), it will be the main type of domain that we shall study. However, we shall show that even for a general domain, there exists a field of fractions or quotient field, thus constructing, in particular, the field of rational numbers. We shall also find the ring analogue of the Chinese Remainder Theorem that we have previously proved for the integers. Finally, we shall define the set theoretic extension of finite induction used in algebra, called Zorn’s Lemma. This axiom, equivalent to the Axiom of Choice that you may have seen, is a powerful tool. [A proof of the equivalence of these two axioms can be found in Appendix A.] We shall then give some simple applications of it that are basic to ring theory and linear algebra.

25. Definitions and Examples

We begin by recalling the definition of a ring, together with specific types of rings.
Definition 25.1. Let $R$ be a set with two binary operations $+ : R \times R \to R$ and $\cdot : R \times R \to R$.

We call $R$ a ring (under $+$ and $\cdot$) if:

1. $(R, +)$ is an additive group (with identity written $0$ or $0_R$).
2. $(R, \cdot)$ is a monoid (with identity written $1$ or $1_R$).
3. $R$ satisfies the distributive laws, i.e., for all $a, b, c \in R$, we have:
   
   (a) $a \cdot (b + c) = a \cdot b + a \cdot c$.
   
   (b) $(b + c) \cdot a = b \cdot a + c \cdot a$.

If (2) is replaced by

(2') $(R, \cdot)$ satisfies associativity (i.e., possibly no 1),

then (following Jacobson), we call $R$ (under $+$ and $\cdot$) a rng. [Take out the condition of having an identity, take out the i from ring.]

A ring satisfying

4. $(R \setminus \{0\}, \cdot)$ is a group

is called a division ring.

As usual if $a, b \in R$, we usually write $ab$ for $a \cdot b$.

Also recall that if $R$ is a ring with $1 \neq 0$, then

$R^\times := \{x \in R \mid \text{There exists } x^{-1} \in R \text{ such that } xx^{-1} = 1 = x^{-1}x\}$

$= \{x \in R \mid \text{There exists } a, b \in R \text{ such that } ax = 1 = xb\}$

is called the group of units of $R$. [Why is it a group?] In particular, a ring is a division ring if and only if $R^\times = R \setminus \{0\}$ (cf. Property 25.4(3c) below).

Definition 25.2. A ring $R$ is called a commutative ring if for all $a, b \in R$, we have

$ab = ba$.

A commutative ring $R$ is called an (integral) domain if

(i) $0 \neq 1$.

(ii) Let $a, b \in R$. If $ab = 0$, then either $a = 0$ or $b = 0$.

A commutative division ring is called a field.

Remark 25.3. Check that a field is a domain.

Rings satisfy the following basic properties (which we leave as an exercise):

Properties 25.4. Let $R$ be a rng and $a, b, c \in R$. Then the following are true:

1. $a \cdot 0 = 0 = 0 \cdot a$. 
(2) \[ a \cdot (-b) = (-a) \cdot b = -(a \cdot b). \]

(3) If \( R \) is a ring, then for all \( a \in R \):
   
   \begin{enumerate}
   
   \item \(-1 \cdot a = -a.\)
   
   \item \((-1) \cdot (-1) = 1.\)
   
   \item \(|R| > 1 \) if and only if \( 1 \neq 0 \).
   
   \end{enumerate}

(4) If \( R \) is a commutative ring with \( 1 \neq 0 \), then \( R \) is a domain if and only if \( R \) satisfies the Cancellation Law: If \( a, b, c \) lie in \( R \), then

\[ a \cdot b = a \cdot c \text{ with } a \neq 0 \text{ implies } b = c. \]

We already know many examples of rings. We recall some of these.

**Examples 25.5.** 1. A **trivial ring** is a ring with one element. By the above, this is the case if and only if \( 1 = 0 \) in \( R \).

2. \( \mathbb{Z} \) is a domain.

3. \( \mathbb{Q}, \mathbb{R}, \mathbb{C} \) are fields.

4. The subset of \( \mathbb{C} \) given by

\[ \mathbb{Z}[-1] = \{ a + b\sqrt{-1} \mid a, b \in \mathbb{Z} \} \]

is a domain under the restriction of the + and \( \cdot \) of \( \mathbb{C} \) called the **Gaussian integers**. If \( z \in \mathbb{Z}[-1] \), there exist unique integers \( a \) and \( b \) satisfying \( z = a + b\sqrt{-1} \). (Cf. the next example.)

5. The subset of \( \mathbb{C} \) given by

\[ \mathbb{Q}[-1] = \{ a + b\sqrt{-1} \mid a, b \in \mathbb{Q} \} \]

is a field under the restriction of the + and \( \cdot \) of \( \mathbb{C} \). [Can you show this?] Note that if \( z \in \mathbb{Q}[-1] \), there exist unique rational numbers \( a \) and \( b \) satisfying \( z = a + b\sqrt{-1} \). Why?

6. \( \mathbb{Z}/n\mathbb{Z} \), with \( n \in \mathbb{Z}^+ \), is a commutative ring. It is a domain if and only if \( n \) is a prime if and only if it is a field.

7. If \( R \) is a ring so is \( \mathbb{M}_n(R) \). If \( R \) is a commutative ring with \( |R| > 1 \), then \( \mathbb{M}_n(R) \) is commutative if and only if \( n = 1 \). (If \( R \) is a rng then so is \( \mathbb{M}_n(R) \) — obvious definition.)

8. If \( R \) is a ring, then \( R[t] \), the set of polynomials with coefficients in \( R \), is a ring under the usual + and \( \cdot \) of polynomials. [What are the usual + and \( \cdot \) of polynomials? What is \( 1_{R[t]} \)?]

9. Let \( R \) be a ring. Define the **ring of (formal) power series** over \( R \) to be the set

\[ R[[t]] := \left\{ \sum_{i=0}^{\infty} a_i t^i \mid a_i \in R \text{ for all } i \right\} \]
with properties and operations:

\[
0_R = 0_R \quad \text{and} \quad 1_R = 1_R, \quad \text{with} \quad t^0 = 1_R
\]

\[
\sum_{i=0}^{\infty} a_i t^i = \sum_{i=0}^{\infty} b_i t^i \quad \text{if and only if} \quad a_i = b_i \quad \text{for all} \quad i
\]

\[
\sum_{i=0}^{\infty} a_i t^i + \sum_{i=0}^{\infty} b_i t^i := \sum_{i=0}^{\infty} (a_i + b_i) t^i
\]

\[
\sum_{i=0}^{\infty} a_i t^i \cdot \sum_{i=0}^{\infty} b_i t^i := \sum_{i=0}^{\infty} c_i t^i \quad \text{with} \quad c_i = \sum_{j=0}^{i} a_j b_{i-j} \quad \text{for all} \quad i.
\]

Since we do not have to worry about convergence, it is easy to define power series (inductively) to satisfy conditions.

Just as for groups, we need to define the maps of interest between rings.

**Definition 25.6.** A map \( \varphi : R \to S \) of rings is called a ring homomorphism if for all \( a, b \in R \), the map \( \varphi \) satisfies:

1. \( \varphi(0_R) = 0_S \).
2. \( \varphi(a + b) = \varphi(a) + \varphi(b) \).
3. \( \varphi(1_R) = 1_S \).
4. \( \varphi(a \cdot b) = \varphi(a) \cdot \varphi(b) \).

I.e., \( \varphi \) is both a group homomorphism \( \varphi : (R, +) \to (S, +) \) and a ‘monoid homomorphism’ \( \varphi : (R, \cdot) \to (S, \cdot) \). Note that we know that (1) is automatic for group homomorphisms. [If we omit (3), we call \( \varphi \) a ring homomorphism.]

A ring homomorphism \( \varphi : R \to S \) is called

i. a ring monomorphism or monic or mono if \( \varphi \) is injective.
ii. a ring epimorphism or epic or epi if \( \varphi \) is surjective.
iii. a ring isomorphism if \( \varphi \) is bijective with inverse a ring homomorphism.
iv. a ring automorphism if \( R = S \) and \( \varphi \) is a ring isomorphism.

Of course, if \( R \) and \( S \) are fields, etc., then \( \varphi \) is called a field homomorphism, etc.

If \( \varphi : R \to S \) is a ring homomorphism, we let

\[
\ker \varphi := \{ x \in R \mid \varphi(x) = 0_S \} \subset R
\]

called the kernel of \( \varphi \), and

\[
\text{im} \varphi := \{ \varphi(x) \mid x \in R \} \subset S
\]

called the image of \( \varphi \).
Remarks 25.7. 1. A ring homomorphism \( \varphi : R \to S \) is a ring isomorphism if and only if \( \varphi \) is bijective. [Proof?]

2. We say two rings \( R \) and \( S \) are isomorphic if there exists an isomorphism \( \varphi : R \to S \). If this is the case, we write \( R \cong S \) and often write the isomorphism as \( \varphi : R \xrightarrow{\cong} S \).

3. A ring homomorphism \( \varphi : R \to S \) is monic if and only if \( \ker \varphi = 0 := \{0\} \).

4. Let \( n \in \mathbb{Z}^+ \), then the canonical map \( \mathbb{Z} \to \mathbb{Z}/n\mathbb{Z} \) is a ring epimorphism with \( \ker \varphi = n\mathbb{Z} = \{ny \mid y \in \mathbb{Z}\} \).

Remark 25.8. Let \( \varphi : R \to S \) be a ring homomorphism. We leave it as an exercise to show that \( \varphi \) is a monomorphism if and only if given any ring homomorphisms \( \psi_1, \psi_2 : T \to R \) with compositions satisfying \( \varphi \circ \psi_1 = \varphi \circ \psi_2 \), then \( \psi_1 = \psi_2 \). (Cf. Exercise 1.12(6).) However, the analogue of Exercise 1.12(7) is false. In particular, suppose that a ring homomorphism \( \varphi : R \to S \) has the property that whenever there exists ring homomorphisms \( \theta_1, \theta_2 : S \to T \) with compositions satisfying \( \theta_1 \circ \varphi = \theta_2 \circ \varphi \) that \( \theta_1 = \theta_2 \). Then this does not necessarily imply that \( \varphi \) is surjective. For example the inclusion map of \( \mathbb{Z} \) into \( \mathbb{Q} \) is a ring homomorphism that satisfies this condition, but is not a surjective ring homomorphism. (Cf. Exercise 26.20(5).) In category theory, being epi is defined by this condition. For this reason, it is no longer common to define a ring epimorphism as a surjective ring homomorphism. However, for consistency of notation, we shall use this older definition, i.e., call a surjective ring homomorphism a ring epimorphism.

We also need the analogues from group theory of subgroups and normal subgroups.

Definition 25.9. Let \( R \) be a ring, \( S \subset R \) a subset with \( S \) a ring. We say \( S \) is a subring of \( R \) if the inclusion map \( \text{inc} : S \to R \) (often written \( \text{inc} : S \hookrightarrow R \)) is a ring monomorphism, so the restriction of the identity map on \( R \) to \( S \) is the inclusion map of \( S \) into \( R \) with \( 1_S = 1_R \) and \( S \) is closed under the ring operations \( + \) and \( \cdot \) on \( R \), meaning that the restriction of the operations on \( R \) to \( S \) have image in \( S \), viz., \( +|_{S \times S} : S \times S \to S \) and \( \cdot|_{S \times S} : S \times S \to S \). [We also have the obvious notion of a subrng of a rng, if we ignore, in this case, the condition on the 1’s, if any.]

Examples 25.10. 1. \( \mathbb{Z} \subset \mathbb{Q} \) and \( \mathbb{Z}, \mathbb{Q}, \mathbb{R} \subset \mathbb{C} \) are subrings.

2. If \( R \) is a domain, e.g., a field, and \( S \subset R \) a subring, then \( S \) is a domain.
3. Let \( \mathbb{Z} \times \mathbb{Z} \) be a ring under componentwise + and \( \cdot \), so \( 1_{\mathbb{Z} \times \mathbb{Z}} = (1, 1) \) and \( \mathbb{Z} \times \mathbb{Z} \) is commutative ring, but not a domain. We also have \( \mathbb{Z} \times 0 := \mathbb{Z} \times \{0\} \subset \mathbb{Z} \times \mathbb{Z} \) is not a subring but is a subrng.

4. Let \( R \) be a ring. The **center** of \( R \) is the set \( Z(R) := \{ a \in R \mid ax = xa \text{ for all } x \in R \} \).

It is a commutative subring of \( R \). For example, the center of a division ring is a field.

**Remark 25.11.** The last example allows us to generalize our terminology. Let \( R \) be a commutative ring, \( A \) a ring and \( \varphi : R \to A \) a ring homomorphism. We call \( A \) an \( R \)-**algebra** via \( \varphi \) if \( \varphi(R) \subset Z(A) \). If \( B \) is another \( R \)-algebra via \( \tau \), then a ring homomorphism \( \psi : A \to B \) is called an \( R \)-**algebra homomorphism** if

\[
\begin{array}{ccc}
A & \xrightarrow{\psi} & B \\
\varphi \downarrow & & \tau \downarrow \\
R & & \\
\end{array}
\]

commutes. Such a \( \tau \) is called an \( R \)-**algebra monomorphism** (respectively, an \( R \)-**algebra epimorphism** \( R \)-**algebra isomorphism**) if it is a ring monomorphism (respectively, epimorphism, isomorphism).

We need the analogue of normal subgroups in ring theory, i.e., those subobjects of a ring that are kernels of ring homomorphisms. If a ring homomorphism \( \varphi : R \to S \) satisfies \( \varphi(1_R) = 0_S \), then \( S = 0 \) as \( \varphi(1_R) = 1_S \). In particular, \( \ker \varphi \) cannot be a subring of \( R \) unless \( \ker \varphi = R \); so in general kernels are not subrings. They are, however, always subrings. Just as in the case when being a subgroup of a group is not sufficient to be the kernel of a group homomorphism, neither is being a subrng sufficient. Rather we need the following objects:

**Definition 25.12.** Let \( R \) be a ring and \( \mathfrak{A} \) a non-empty subset of \( R \). We call \( \mathfrak{A} \) a \((2\text{-sided})\) **ideal** of \( R \) if \( (\mathfrak{A}, +) \subset (R, +) \) is a subgroup and for all \( r \) in \( R \) and \( a \) in \( \mathfrak{A} \), we have both \( ar \) and \( ra \) lie in \( \mathfrak{A} \), i.e., \( \cdot : R \times \mathfrak{A} \to \mathfrak{A} \) and \( \cdot : \mathfrak{A} \times R \to \mathfrak{A} \). So an ideal of \( R \) is a subrng of \( R \) closed under multiplication by \( R \) on both the right and left. If \( \mathfrak{A} \) only satisfies being a subrng and closed under multiplication by \( R \) on the left (respectively, on the right), we call \( \mathfrak{A} \) a **left ideal** (respectively, **right ideal**). Of course, if \( R \) is a commutative ring, the notions of \((2\text{-sided})\) ideal, right ideal, and left ideal coincide, making life much simpler.

**Examples 25.13.** Let \( R \) be a nontrivial ring.
25. DEFINITIONS AND EXAMPLES

1. \(0 = \{0\}\) is an ideal of \(R\) called the \textit{trivial ideal} and \(R\) is an ideal of \(R\) called the \textit{unit ideal}. If these are the only ideals of \(R\), then we call \(R\) a \textit{simple ring}. (Cf. simple groups.)

2. Let \(a \in R\). Define\[
(a) = RaR := \{\sum_{i=1}^{n} x_i ay_i \mid n \in \mathbb{Z}^+, x_i, y_i \in R, 1 \leq i \leq n\},
\]
an ideal of \(R\) called the \textit{principal ideal generated} by \(a\). It is the smallest ideal in \(R\) containing \(a\). [Cf. this with a cyclic subgroup of a group.] If \(R\) is commutative, this is much simpler, as (using the distributive laws), it is just \((a) = Ra := \{ra \mid r \in R\}\).

3. Let \(a_1, \ldots, a_n \in R\). The smallest ideal \textit{generated} by \(a_1, \ldots, a_n\) is \[(a_1, \ldots, a_n) = Ra_1 R + \cdots + Ra_n R := \{\sum_{i=1}^{n} \sum_{j=1}^{m_i} x_{ij} a_i y_{ij} \mid x_{ij}, y_{ij} \in R \text{ all } i, j \text{ some } n, m_i \in \mathbb{Z}^+\}.
\]

If \(R\) is commutative, this is simply \[(a_1, \ldots, a_n) = Ra_1 + \cdots + Ra_n := \{\sum_{i=1}^{n} r_i a_i \mid r_i \in R, 1 \leq i \leq n\},
\]
i.e., it is all \(R\)-linear combinations of \(a_1, \ldots, a_n\). [Cf. this with the span of finitely many vectors in a vector space.]

4. Let \(I\) be an (indexing) set, \(A = \{a_i \mid i \in I\}\). Then the smallest ideal in \(R\) generated by \(A\) is\[
\langle A \rangle = (a_i)_I := \{a \in R \mid \text{There exist } a_{i_1}, \ldots, a_{i_n} \in A, \text{ some } n \text{ and } i_j \in I, \text{ satisfying } a \in (a_{i_1}, \ldots, a_{i_n})\}.
\]
[So \(\langle A \rangle\) is the union of all the ideals generated by \(a_{i_1}, \ldots, a_{i_n}\) in \(A\) with \(i_1, \ldots, i_n \in I\), for some \(n\).]

5. Let \(\mathfrak{A}\) and \(\mathfrak{B}\) be ideals in \(R\). Then\[
\mathfrak{A} + \mathfrak{B} := (\mathfrak{A} + \mathfrak{B}) = \{a + b \mid a \in \mathfrak{A} \text{ and } b \in \mathfrak{B}\}
\]
\[
\mathfrak{A} \cap \mathfrak{B} := \{x \mid x \in \mathfrak{A} \text{ and } x \in \mathfrak{B}\}
\]
\[
\mathfrak{A} \mathfrak{B} := (\{ab \mid a \in \mathfrak{A}, b \in \mathfrak{B}\})
\]
\[
= \{\sum_{i=1}^{n} a_i b_i \mid a_i \in \mathfrak{A}, b_i \in \mathfrak{B}, 1 \leq i \leq n, \text{ some } n\}
\]
are ideals in $R$. Moreover, we have

$$\mathfrak{A}\mathfrak{B} \subset \mathfrak{A} \cap \mathfrak{B} \subset \mathfrak{A} + \mathfrak{B}$$

(and they are usually all different).

Note the formal difference between the definition of $\mathfrak{A}\mathfrak{B}$ for rings and $\mathfrak{H}\mathfrak{K}$, with $\mathfrak{H}$, $\mathfrak{K}$ subgroups of $G$, although this difference vanishes if $\mathfrak{H}$ or $\mathfrak{K}$ is a normal subgroup (the analogue of an ideal). So $\mathfrak{A}\mathfrak{B}$ is the ideal in $R$ generated by $ab$ with $a \in \mathfrak{A}$, $b \in \mathfrak{B}$. More generally, we can recursively define the product of finitely many ideals. [Note that infinite products do not make sense in a ring.]

Note $\mathfrak{A} \cup \mathfrak{B}$ is usually not an ideal (cf. groups and vector spaces). Can you give a condition that will guarantee that it is?

6. If $\varphi : R \to S$ is a ring homomorphism, then $\ker \varphi$ is an ideal of $R$. But in general $\text{im} \varphi$ is not an ideal of $S$. Indeed, it is an ideal if and only if $\varphi$ is surjective. It is, however, a subring of $S$.

7. If $\mathfrak{A}$ is an ideal of $R$ and $\mathfrak{A} \cap R^\times \neq \emptyset$, then $\mathfrak{A} = R$, the unit ideal (and conversely). Indeed, if $x \in \mathfrak{A} \cap R^\times$ and $r \in R$, then $r = r1 = rx^{-1}x \in \mathfrak{A}$.

We now give some basic examples.

**Specific Examples 25.14.** Let $R$ be a ring.

1. If $R$ is a division ring, then $R$ is simple. Indeed, if $0 < \mathfrak{A} \subset R$ is an ideal, then $\mathfrak{A} \cap R^\times \neq \emptyset$.

2. Suppose $R$ is not the trivial ring. If $R$ is commutative, then $R$ is simple if and only if $R$ is a field. Indeed if $R$ is simple, $0 \neq a \in R$, then $(a) = R$; so there exists an $r \in R$ satisfying $1 = ar = ra$ showing $a$ is a unit. The converse is the previous example.

3. $\mathfrak{M}_n(\mathbb{R})$, for $n > 1$, is simple and not a division ring. In fact, if $R$ is simple, so is $\mathfrak{M}_n(R)$ which is never a division ring if $n > 1$. [We leave this as an exercise.]

4. Let $R = \mathbb{Z}$ and $0 < \mathfrak{A} \subset \mathbb{Z}$ an ideal. Then $(\mathfrak{A}, +) \subset (\mathbb{Z}, +)$ is a subgroup, so $\mathfrak{A} = \mathbb{Z}n$ for some $n \in \mathbb{Z}^+$ as an additive group. Clearly, $(n) = \mathbb{Z}n \subset \mathbb{Z}$ is an ideal. Therefore, every ideal in $\mathbb{Z}$ is of the form $(n)$ for some $n \geq 0$ in $\mathbb{Z}$. In particular, every ideal in $\mathbb{Z}$ is principal, where an ideal $\mathfrak{A}$ is called principal if there exists an $a \in \mathfrak{A}$ such that $\mathfrak{A} = (a)$.

5. A domain $R$ is called a principal ideal domain or a PID if every ideal in $R$ is principal, e.g., $\mathbb{Z}$ is a PID.

6. We give further examples of PIDs (without proof):
(i) Any field.
(ii) $F[t]$, if $F$ is a field. [Can you show this?]
(iii) Let $\mathbb{Z}[\theta] := \{a + b\theta \mid a, b \in \mathbb{Z}\} \subset \mathbb{C}$ with
$$\theta = \sqrt{-1}, \sqrt{-2}, \frac{-1 - \sqrt{-3}}{2}, \frac{-1 - \sqrt{-7}}{2}, \frac{-1 - \sqrt{-11}}{2},$$
$$\frac{-1 - \sqrt{-19}}{2}, \frac{-1 - \sqrt{-43}}{2}, \frac{-1 - \sqrt{-67}}{2}, \frac{-1 - \sqrt{-163}}{2}.$$ 
Under the $+$ and $\cdot$ in $\mathbb{C}$, these are rings and each is a PID.
Gauss had conjectured that these were the only PID’s of this form (with negative square roots). It was proved by Harold Stark.

7. $\mathbb{Z}[\sqrt{-5}] := \{a + b\sqrt{-5} \mid a, b \in \mathbb{Z}\} \subset \mathbb{C}$ is a domain but not a PID.
8. $\mathbb{Z}[t]$ is a domain but not a PID.

We begin our study by generalizing concepts that we investigated when studying the integers.

**Definition 25.15.** Let $R$ be a commutative ring, $a, b$ in $R$ with $a \neq 0$. We say that $a$ divides $b$ and write $a \mid b$, if there exists an $x$ in $R$ satisfying $b = ax$.

**Properties 25.16.** Let $R$ be a commutative ring and $a, b, c$ elements in $R$.

1. If $a \mid b$ and $a \mid c$, then $a \mid bx + cy$ for all $x, y$ in $R$.
2. If $R$ is a domain, then $a \mid b$ and $b \mid a$ if and only if there exists a unit $u$ in $R$ satisfying $a = ub$.

Note the second property generalizes the corresponding result for the integers, as $\pm 1$ are the only units in $\mathbb{Z}$.

**Proof.** The proof of the first property is the same as for the integers, so we need only prove the second. If $a = ub$ for a unit $u$ then $b = u^{-1}a$. Conversely, if $a \mid b$ and $b \mid a$ then there exist $x, y \in R$ satisfying $b = ax$ and $a = by$. As $a \neq 0$ in the domain $R$ and $a = by = axy$, we have $1 = xy$ by the Cancellation Law, so $x, y \in R^\times$. 

Property (2) motivates the following definition: If $R$ is a domain, $a$ and $b$ nonzero elements in $R$, we say that $a$ is an associate of $b$ and write $a \approx b$ if there exists a unit $u$ in $R$ such that $a = ub$. Thus $a$ is an associate of $b$ if and only if $a \mid b$ and $b \mid a$. [Check that $\approx$ is an
equivalence relation.] For example, if $R = \mathbb{Z}$, then the only associates of $n$ in $\mathbb{Z}$ are $\pm n$. (We also let 0 be the associate of 0.)

**Remarks 25.17.** Let $R$ be a commutative ring. Then the relationship between division and principal ideal is given by the following: Let $a$, $b$, and $c$ be nonzero elements in $R$.

1. $a \in (b)$ if and only if $b \mid a$ if and only if $(a) \subset (b)$.
2. $(a) = (b)$ if and only if $a \mid b$ and $b \mid a$.
3. $a \mid b$ and $a \mid c$ if and only if $(b, c) \subset (a)$.
4. $(b, c) = (a)$ if and only if $a \mid b$ and $a \mid c$ and there exist $x$ and $y$ in $R$ satisfying $a = bx + cy$.
5. Suppose that $R$ is a domain. Then
   (a) $a \approx b$ if and only if $(a) = (b)$.
   (b) If $a = rb$ in $R$, then $a \approx b$ if and only if $r \in R^\times$. In particular, if $a \not\approx b$ and $a = rb$ in $R$, then $(a) < (b)$. Indeed, if $ub = a = rb$ with $u \in R^\times$, then $u = r$ as $R$ is a domain.

If $R = \mathbb{Z}$ and $p \neq \pm 1$ is an integer, we know that $p$ is a prime if and only if $p \mid ab$ implies $p \mid a$ or $p \mid b$. This motivates the most important definition concerning ideals.

**Definition 25.18.** Let $R$ be a commutative ring and $\mathfrak{A} < R$ an ideal. We call $\mathfrak{A}$ a **prime ideal** in $R$ if

$$ab \in \mathfrak{A} \implies a \in \mathfrak{A} \text{ or } b \in \mathfrak{A}.$$ 

It is called a **maximal ideal** in $R$ if

$$\mathfrak{A} < \mathfrak{B} \subset R \text{ is an ideal, then } \mathfrak{B} = R.$$ 

[Note this definition of maximal ideal (even left maximal ideal — definition?) makes sense in any ring, commutative or not.]

**Examples 25.19.** Let $R$ be a nontrivial commutative ring.

1. $R$ is a domain if and only if $0 = (0)$ is a prime ideal in $R$:
   We have 0 is a prime ideal if and only if $ab \in (0)$ implies $a \in (0)$ or $b \in (0)$ if and only if $ab = 0$ implies $a = 0$ or $b = 0$.
2. Every maximal ideal $\mathfrak{m} < R$ is a prime ideal:
   Let $ab \in \mathfrak{m}$ with $a \not\in \mathfrak{m}$. Then $\mathfrak{m} < \mathfrak{m} + Ra \subset R$ is an ideal, hence $R = \mathfrak{m} + Ra$. Therefore, there exist $m \in \mathfrak{m}$ and $r \in R$ satisfying $1 = m + ra$. Hence $b = bm + rab \in \mathfrak{m} + m \subset \mathfrak{m}$.
3. Suppose that $R$ is a PID. Then every nonzero prime ideal in $R$ is maximal:
   [Warning: This is not true in general, e.g., in $\mathbb{Z}[t]$ — why?]
Let $0 < p < R$ be a prime ideal and suppose that $p < A \subset R$ is an ideal. By the definition of a PID, there exist $p \in p$ and $a \in A$ satisfying $p = (p)$ and $A = (a)$. Therefore, we have $p = ra$ for some $r \in R$ and $a \notin (p)$. As $ra = p \in p$, we have $r \in p$. Write $r = sp$, with $s \in R$. Then $p = ra = ps a$ in the domain $R$, so $1 = sa$. Hence $a \in R^\times$ and $A = R$. This shows that $p$ is maximal.

4. The prime ideals in $\mathbb{Z}$ are:

- $0 = (0)$ prime, not maximal
- $p\mathbb{Z} = (p)$ a prime, maximal.

We want to view prime ideals as an ideal theoretical property rather than a property about elements in the ideal. This is easily done and is the direct analogue of Euclid’s Lemma for ideals.

**Lemma 25.20.** Let $R$ be a commutative ring and $p < R$ an ideal. Then $p$ is a prime ideal if and only if for all ideals $A$ and $B$ in $R$ satisfying $AB \subset p$, either $A \subset p$ or $B \subset p$. Moreover, if $A < R$ is not a prime ideal, then there exists ideals $A_i < R$, $i = 1, 2$ satisfying $A_1A_2 \subset A$.

**Proof.** ($\Leftarrow$): If $xy \in p$ with $x, y \in R$ then $(x)(y) \subset p$, so either $x \in (x) \subset p$ or $y \in (y) \subset p$.

($\Rightarrow$): Suppose that the ideal $A \not\subset p$ but $AB \subset p$ with $B$ an ideal. Then there exists an element $a \in A \setminus p$. As $ab \in p$ for all $b \in B$, it follows that $B \subset p$.

As for the last statement, suppose that $A < R$ is not a prime ideal. Then there exist ideals $B_i \not\subset A$, $i = 1, 2$ but $B_1B_2 \subset A$. Let $B_i' = B_i + A > A$, then $B_1'B_2' \subset B_1B_2 + A \subset A$. □

**Exercises 25.21.**

1. Prove if $R$ is a domain so is the ring of formal power series $R[[t]]$.

2. Let $R$ be a commutative ring. Show that if $f = 1 + \sum_{i=1}^{\infty} a_i$ is a formal power series in $R[[t]]$, then one can determine $b_1, \ldots, b_n, \ldots$ such that $g = 1 + \sum_{i=1}^{\infty} b_i$ is the multiplicative inverse of $f$ in $R[[t]]$. In particular,

$$R[[t]]^\times = \{ a_0 + \sum_{i=1}^{\infty} a_i t^i \in R[[t]] \mid a_0 \in R^\times \}.$$  

3. Let $R$ be a commutative ring and $A = R[[t]]$, the power series over $R$. Let $\mathfrak{p}$ be a prime ideal in $A$ not containing $t$ and $\varphi : A \to R$ the homomorphism induced by $t \mapsto 0$. Suppose that $\varphi(\mathfrak{p}) = (a_1, \ldots, a_n)$. Let $f_i = a_i +$ higher terms in $t$ for $i = 1, \ldots, n$. Prove
if \( g \in P \), then \( g = f_1 h_1 + \cdots + f_n h_n \) for some elements \( h_1, \ldots, h_n \) in \( A \).

4. Let \( A \) be an additive group and let \( \text{End}(A) \) denote the set of group homomorphisms of \( A \) to \( A \). Prove that \( \text{End}(A) \) is a ring under addition and composition of elements in \( \text{End}(A) \). Also show the unit group of \( \text{End}(A) \) is the group of group automorphisms of \( A \).

5. Let \( F \) be a field. Show that \( \mathbb{M}_n(F) \) is a simple ring. If \( n > 1 \) then \( \mathbb{M}_n(F) \) is not a division ring.

6. Show that a ring homomorphism \( \phi : R \to S \) is a monomorphism if and only if given any ring homomorphisms \( \psi_1, \psi_2 : T \to R \) with compositions satisfying \( \phi \circ \psi_1 = \phi \circ \psi_2 \), then \( \psi_1 = \psi_2 \).

7. Show if \( R \) is a commutative ring, then, for any ideal \( \mathfrak{B} \) in \( \mathbb{M}_n(R) \), there exists an ideal \( \mathfrak{A} \) in \( R \) satisfying \( \mathfrak{B} = \mathbb{M}_n(\mathfrak{A}) \) (obvious definition). In particular, if \( R \) is simple, so is \( \mathbb{M}_n(R) \). Show that \( \mathbb{M}_n(R) \) is never a division ring if \( n > 1 \).

8. Let \( R \) be a commutative ring. Prove the Binomial Theorem: Let \( a \) and \( b \) be elements of \( R \) and \( n \) a positive integer. Then

\[
(a + b)^n = \sum_{i=0}^{n} \binom{n}{i} a^i b^{n-i}.
\]

9. If \( R \) is a ring satisfying \( x^2 = x \) for all \( x \) in \( R \), then \( R \) is commutative.

10. If \( R \) is a rng satisfying \( x^3 = x \) for all \( x \) in \( R \), then \( R \) is commutative.

11. Let \( R \) be a commutative ring and \( \mathfrak{A} \) be an ideal in \( R \) satisfying

\[
\mathfrak{A} = m_1 \cdots m_r = n_1 \cdots n_s
\]

with all the \( m_i \) distinct maximal ideals and all the \( n_j \) distinct maximal ideals. Show that \( r = s \) and there exists a \( \sigma \in S_r \) satisfying \( m_i = n_{\sigma(i)} \) for all \( i \).

12. Let \( R \) be a commutative ring and \( \mathfrak{A} \) an ideal of \( R \). Suppose that every element in \( R \setminus \mathfrak{A} \) is a unit of \( R \). Show that \( \mathfrak{A} \) is a maximal ideal of \( R \) and that, moreover, it is the only maximal ideal of \( R \).

13. Let \( R \) be the set of all continuous functions \( f : [0, 1] \to \mathbb{R} \). Then \( R \) is a commutative ring under + and \( \cdot \) of functions. Show that any maximal ideal of \( R \) has the form \( \{ f \in R \mid f(a) = 0 \} \) for some fixed \( a \) in \( [0, 1] \).

14. Let \( \varphi : R \to S \) be a ring homomorphism of commutative rings. Show that if \( \mathfrak{B} \) is an ideal (respectively, prime ideal) of \( S \) then \( \varphi^{-1}(\mathfrak{B}) \) is an ideal (respectively, prime ideal) of \( R \). Give an example with \( \mathfrak{B} \) a maximal ideal of \( S \) but \( \varphi^{-1}(\mathfrak{B}) \) not a maximal ideal of \( R \).
15. Show \( F[t] \), with \( F \) a field, is a PID.

16. Show that \( \mathbb{Z}[t] \), the ring of polynomials with coefficients in \( \mathbb{Z} \), is not a PID. Show this by finding a non-principal maximal ideal in \( \mathbb{Z}[t] \). Also find a nonzero prime ideal in \( \mathbb{Z}[t] \) that is not maximal and prove it is such.

17. Let \( \mathfrak{A}_1, \ldots, \mathfrak{A}_n \) be ideals in \( R \), at least \( n - 2 \) of which are prime. Let \( S \subset R \) be a subrng (it does not have to have a 1) contained in \( \mathfrak{A}_1 \cup \cdots \cup \mathfrak{A}_n \). Then one of the \( \mathfrak{A}_j \)'s contains \( S \). In particular, if \( p_1, \ldots, p_n \) are prime ideals in \( R \) and \( \mathfrak{B} \) is an ideal properly contained in \( S \) satisfying \( S \setminus \mathfrak{B} \subset p_1 \cup \cdots \cup p_n \), then \( S \) lies in one of the \( p_i \)'s.

26. Factor Rings and Rings of Quotients

We begin by generalizing the construction of the integers modulo \( n \).

**Definition 26.1.** Let \( R \) be a ring, \( \mathfrak{A} \) an ideal in \( R \). If \( a \) and \( b \) are elements in \( R \), write \( a \equiv b \mod \mathfrak{A} \) if \( a - b \) is an element of \( \mathfrak{A} \).

A proof analogous proof to that for the integers modulo \( n \), establishes the following:

**Proposition 26.2.** Let \( R \) be a ring and \( \mathfrak{A} \) an ideal in \( R \). Then \( \equiv \mod \mathfrak{A} \) is an equivalence relation. Suppose that the elements \( a, a', b, b' \) in \( R \) satisfy:

\[
\begin{align*}
a \equiv a' \mod \mathfrak{A} \quad \text{and} \quad b \equiv b' \mod \mathfrak{A},
\end{align*}
\]

then

\[
\begin{align*}
a + b & \equiv a' + b' \mod \mathfrak{A}, \\
a \cdot b & \equiv a' \cdot b' \mod \mathfrak{A}.
\end{align*}
\]

Let \( \overline{R} = R/\mathfrak{A} := R/(\equiv \mod \mathfrak{A}) \). If \( a \in R \), set

\[
\overline{a} = \{ b \in R \mid b \equiv a \mod \mathfrak{A} \} = a + \mathfrak{A}.
\]

If

\[
\overline{\cdot} : R \to R/\mathfrak{A} \quad \text{given by} \quad a \mapsto \overline{a}
\]

denotes the canonical surjection and \( a, b \) lie in \( R \), define

\[
\begin{align*}
(i) \quad \overline{a + b} & := \overline{a + b} \\
(ii) \quad \overline{a \cdot b} & := \overline{a \cdot b}.
\end{align*}
\]
Then under these operations, $\overline{R}$ is a ring with $0_{\overline{R}} = \mathfrak{A}$ and $1_{\overline{R}} = 1_R + \mathfrak{A}$ with $\overline{R}$ commutative if $R$ is. The canonical surjection is an epimorphism with kernel $\mathfrak{A}$. We call $\overline{R}$ the factor or quotient ring of $R$ by $\mathfrak{A}$.

Proofs analogous to those in Group Theory show:

**Theorem 26.3.** (First Isomorphism Theorem) Let $\varphi : R \to S$ be a ring homomorphism. Then we have a commutative diagram of rings and ring homomorphisms

$$
\begin{array}{ccc}
R & \xrightarrow{\varphi} & S \\
\downarrow \varphi & & \downarrow \text{inc} \\
R/\ker \varphi & \xrightarrow{\overline{\varphi}} & \text{im } \varphi
\end{array}
$$

with $\overline{\varphi}$ a ring isomorphism between $R/\ker \varphi$ and $\text{im } \varphi$, and $\text{inc}$ a ring monomorphism.

and

**Theorem 26.4.** (Correspondence Principle) Let $\varphi : R \to S$ be a ring epimorphism. Then

$$\{ \mathfrak{B} \mid \mathfrak{B} \text{ an ideal in } R \text{ with } \ker \varphi \subset \mathfrak{B} \} \longrightarrow \{ \mathfrak{C} \mid \mathfrak{C} \text{ an ideal in } S \}$$
given by $\mathfrak{B} \mapsto \varphi(\mathfrak{B})$ is an order preserving bijection.

We give some applications.

**Proposition 26.5.** Let $R$ be a nontrivial commutative ring and $\mathfrak{A}$ an ideal of $R$. Then

1. $\mathfrak{A}$ is a prime ideal of $R$ if and only if $R/\mathfrak{A}$ is a domain.
2. $\mathfrak{A}$ is a maximal ideal of $R$ if and only if $R/\mathfrak{A}$ is a field.

**Proof.** Let $\overline{\cdot} : R \to R/\mathfrak{A}$ be the canonical epimorphism, so $\ker \overline{\cdot} = \mathfrak{A}$.

1: $\overline{R}$ is a domain if and only if $(\overline{0})$ is a prime ideal in $\overline{R}$ if and only if whenever $a, b \in R$ satisfy $\overline{a} \cdot \overline{b} = \overline{0}$, then $\overline{a} \in (\overline{0})$ or $\overline{b} \in (\overline{0})$ if and only if whenever $a, b \in R$ satisfy $a \cdot b \in \mathfrak{A}$, then $a \in \mathfrak{A}$ or $b \in \mathfrak{A}$ if and only if $\mathfrak{A}$ is a prime ideal.

2: $(\overline{0}) \neq \overline{R}$ is a field if and only if $\overline{R}$ is simple if and only if $(\overline{0})$ and $\overline{R}$ are the only ideals of $\overline{R}$ if and only if $\{ \mathfrak{B} \mid \mathfrak{A} \subset \mathfrak{B} \subset R \text{ is an ideal} \} = \{ \mathfrak{A}, R \}$ if and only if $\mathfrak{A} < R$ is maximal by the Correspondence Principle. \qed
Question Is (2) true if the ring $R$ is not necessarily commutative and we replace field by division ring?

Compare the next application with the Third Isomorphism Theorem of Groups.

**Proposition 26.6.** Let $R$ be a ring and $\mathfrak{A} \subset \mathfrak{B} \subset R$ ideals. Then

$$\mathfrak{B}/\mathfrak{A} := \{ b + \mathfrak{A} \mid b \in \mathfrak{B} \} \subset R/\mathfrak{A}$$

is an ideal in $R/\mathfrak{A}$ and

$$R/\mathfrak{B} \cong (R/\mathfrak{A})/(\mathfrak{B}/\mathfrak{A}).$$

**Proof.** It is easy to check that $R/\mathfrak{A} \to R/\mathfrak{B}$ given by $r + \mathfrak{A} \mapsto r + \mathfrak{B}$ is a well-defined epimorphism with kernel $\mathfrak{B}/\mathfrak{A}$. □

**Proposition 26.7.** Let $\varphi : R \to S$ be a ring epimorphism of commutative rings with kernel $\mathfrak{A}$. Then the map

$$\{ p \mid p < R \text{ a prime ideal with } \mathfrak{A} \subset p \} \to \{ \mathfrak{P} \mid \mathfrak{P} < S \text{ a prime ideal} \}$$

given by $p \mapsto \varphi(p)$ is a bijection.

**Proof.** We may assume that $S$ is not the trivial ring for if not then both sets of prime ideals would be empty. By the First Isomorphism Theorem, we may assume that $S = R/\mathfrak{A}$ and $\varphi$ is the canonical epimorphism $- : R \to R/\mathfrak{A}$. Let $\mathfrak{A} \subset \mathfrak{B} \subset R$ be an ideal. Then $\mathfrak{B}$ is a prime ideal in $R$ if and only if $R/\mathfrak{B}$ is a domain if and only if $(R/\mathfrak{A})/(\mathfrak{B}/\mathfrak{A})$ is a domain if and only if $\mathfrak{B}/\mathfrak{A}$ is a prime ideal in $R/\mathfrak{A}$. □

**Warning 26.8.** If $\varphi : R \to S$ is a ring homomorphism of rings with $p$ a prime ideal in $R$, it does not follow that the ideal generated by image of $p$, $S\varphi(p)$, is a prime ideal in $S$. For example, the inclusion $\mathbb{Z} \subset \mathbb{Q}$ is a ring homomorphism and (2) is a prime ideal in $\mathbb{Z}$, but its image $(2)$ is the unit ideal in $\mathbb{Q}$.

**Construction 26.9.** Let $R$ be a nontrivial ring. As any ring homomorphism must take the multiplicative identity to the multiplicative identity, there exists a unique ring homomorphism from $\mathbb{Z}$ to $R$. It is given by

$$\iota : \mathbb{Z} \to R \text{ given by } m \mapsto m1_R = \begin{cases} 1_R + \cdots + 1_R & \text{if } m \geq 0 \\ -1_R + \cdots + -1_R & \text{if } m < 0. \end{cases}$$

Since $\mathbb{Z}$ is a PID, there exists a unique integer $n \geq 0$ such that $\ker \iota = (n) = (-n)$. We call $n$ the characteristic of the ring $R$ and write $\text{char}(R)$ (or just $\text{char } R$) for this integer. We shall only be interested in the characteristic of commutative rings.
Remarks 26.10. Let \( R \) be a commutative ring and \( \iota : \mathbb{Z} \rightarrow R \) the unique ring homomorphism.

1. \( \text{char}(R) = 0 \) if and only if \( m1_R \neq 0 \) for all nonzero integers \( m \) if and only if \( \iota \) is monic.

2. If \( \text{char}(R) \neq 0 \) then, by the Division Algorithm, \( \text{char}(R) \) is the least positive integer \( n \) such that \( n1_R = 0 \).

3. If \( \varphi : \mathbb{Z}/n\mathbb{Z} \rightarrow R \) and \( \psi : \mathbb{Z}/m\mathbb{Z} \rightarrow R \) are ring monomorphisms with \( n, m \) positive integers, then \( n = \text{char}(R) = m \) and \( \varphi = \psi \) are induced by \( \iota \).

Examples 26.11. 1. \( \mathbb{Z} \) is a domain, and \( \mathbb{Q}, \mathbb{R}, \) and \( \mathbb{C} \) are fields of characteristic zero.

2. If \( n \in \mathbb{Z}^+ \) then \( \text{char}(\mathbb{Z}/n\mathbb{Z}) = n \).

3. If \( \varphi : R \rightarrow S \) is a ring monomorphism, then \( \text{char}(R) = \text{char}(S) \).

4. The canonical epimorphism \( - : \mathbb{Z} \rightarrow \mathbb{Z}/n\mathbb{Z} \) shows that, in general, the characteristic of a ring is not preserved under a ring homomorphism.

Proposition 26.12. Let \( R \) be a domain. Then either

1. \( \text{char}(R) = 0 \) and there exists a ring monomorphism \( \mathbb{Z} \rightarrow R \) or
2. \( \text{char}(R) = p, p \) a prime, and there exists a ring monomorphism \( \mathbb{Z}/p\mathbb{Z} \rightarrow R \).

Proof. Let \( \iota : \mathbb{Z} \rightarrow R \) be the unique ring homomorphism. Then \( \ker \iota = (n) \) for some non-negative integer \( n \), so \( \text{char}(R) = n \). As \( R \) is a domain so is \( \text{im} \iota \cong \mathbb{Z}/n\mathbb{Z} \). Therefore, \( n = 0 \) or \( n \) is a prime and the induced map \( \tilde{\iota} : \mathbb{Z}/n\mathbb{Z} \rightarrow R \) is a monomorphism. \( \square \)

The proposition says that if \( R \) is a domain, we may view \( \mathbb{Z} \subset R \) if \( \text{char}(R) = 0 \) and \( \mathbb{Z}/p\mathbb{Z} \subset R \) if \( \text{char}(R) = p \).

We have now generalized the construction of \( \mathbb{Z}/n\mathbb{Z} \) to \( R/\mathfrak{A} \), where \( \mathfrak{A} \) is an ideal in \( R \). Our next goal is to construct \( \mathbb{Q} \) from \( \mathbb{Z} \). Our construction will work for any domain.

Construction 26.13. Let \( R \) be a domain. We wish to construct fractions from \( R \), i.e., if \( a, b, c, d \in R \) with \( b \neq 0 \) and \( d \neq 0 \), we want to have the following:
\( \frac{a}{b} + \frac{c}{d} = \frac{ad + bc}{bd}. \) [Note \( bd \neq 0 \) in the domain \( R \).]

\( \frac{a}{b} \cdot \frac{c}{d} = \frac{ac}{bd} \).

and we know that we must have

\( \frac{a}{b} = \frac{c}{d} \) if and only if \( ad = bc \) in \( R \).

We do this construction as follows: Let

\[ S = \{ (a, b) \mid a \in R, 0 \neq b \in R \} = R \times (R \setminus \{0\}) \]

and

\( (a, b) \sim (c, d) \) in \( S \) if \( ad = bc \).

Check that \( \sim \) is an equivalence relation.

Let

\[ \frac{a}{b} := [(a, b)]_\sim, \] the equivalence class of \((a, b)\) under \( \sim \) and

\[ K := S/\sim = \{ \frac{a}{b} \mid (a, b) \in S \} = \{ \frac{a}{b} \mid a, b \in R, b \neq 0 \}. \]

Then \( K \) satisfies \((iii)\) above. Define \( + : K \times K \to K \) and \( \cdot : K \times K \to K \) by \((i)\) and \((ii)\) above respectively.

**Claim.** \( + \) and \( \cdot \) are well-defined:

Suppose that \((a, b) \sim (a', b')\) and \((c, d) \sim (c', d')\) in \( S \). We must show

\[ (ad + bc)b'd' = (a'd' + b'c')(bd) \]

knowing that

\[ ab' = a'b \text{ and } cd' = c'd, \]

which is easily checked. Thus \( K \) has \( + \) and \( \cdot \). It is routine to check that \( K \) is a commutative ring with \( 0_K = \frac{0}{b} \), and \( 1_K = \frac{b}{b} \) with \( 0 \neq b \in R \).

If \( \frac{a}{b} \neq 0_K \) then \( a \neq 0 \) in \( R \), hence \( \frac{a}{b} \) is defined and is the inverse of \( \frac{b}{a} \), so \( K \) is a field called the *field of quotients* or the *quotient field* of \( R \).

We write \( qf(R) \) for \( K \).

**Example.** \( \mathbb{Q} = qf(\mathbb{Z}) \).

This construction of \( qf(R) \) is essentially unique as \( qf(R) \) has the following *universal property*:
Theorem 26.14. (The Universal Property of the Field of Quotients)

Let \( R \) be a domain. Then there exists a field \( K \), and a ring monomorphism \( i : R \to K \) satisfying the following Universal Property: If \( F \) is a field, then for any ring monomorphism \( \varphi : R \to F \), there exists a unique ring monomorphism \( \psi : K \to F \) such that the diagram

\[
\begin{array}{ccc}
R & \xrightarrow{i} & K \\
\downarrow{\varphi} & & \downarrow{\psi} \\
F & & 
\end{array}
\]

commutes.

In particular, if \( K' \) is a field and \( j : R \to K' \) a ring monomorphism also satisfying the universal property (\( \dagger \)), then there exists a unique (field) isomorphism \( \sigma : K \to K' \) such that the diagram

\[
\begin{array}{ccc}
R & \xrightarrow{i} & K \\
\downarrow{j} & & \downarrow{\sigma} \\
K' & & 
\end{array}
\]

commutes.

Proof. Let \( K = qf(R) \), then \( i : R \to K \) given by \( r \mapsto \frac{r}{1} \) is a ring homomorphism. Let \( b \) be a nonzero element in \( R \). Then \( \frac{r}{1} = \frac{0}{b} \) in \( K \) if and only if \( rb = 1 \cdot 0 = 0 \) in the domain \( R \) if and only if \( r = 0 \), so this homomorphism is injective.

Now suppose that we are given a ring monomorphism \( \varphi : R \to F \) with \( F \) a field. Define

\[
\psi : K \to F \text{ by } \psi(\frac{a}{b}) = \varphi(a)\varphi(b)^{-1}.
\]

This makes sense, as \( \varphi(b) \neq 0 \), since \( \varphi \) is monic and \( F \) is a field.

Claim 1. \( \psi \) is well-defined and monic.

\[
a = \frac{a'}{b'} \text{ in } K \text{ if and only if } ab' = a'b \text{ in } R \text{ if and only if } \varphi(a)\varphi(b') = \varphi(a')\varphi(b) \text{ in } F \text{ as } \varphi \text{ is monic if and only if } \varphi(a)\varphi(b)^{-1} = \varphi(a')\varphi(b')^{-1} \text{ in the field } F \text{ if and only if } \psi\left(\frac{a}{b}\right) = \psi\left(\frac{a'}{b'}\right) \text{ in } F. \]

It is easily seen that \( \psi \) is a homomorphism. The Claim follows.

Claim 2. \( \psi \) is unique.

If \( \psi' \) is another monomorphism making the diagram (\( \dagger \)) commute, then

\[
\psi'(\frac{a}{1}) = \varphi(a) = \psi\left(\frac{a}{1_K}\right) \text{ for all } a \in R
\]
in the field \( F \). In particular, if \( b \neq 0 \) in \( R \), we have
\[
\psi'\left(\frac{b}{1}\right)\psi'\left(\frac{1}{b}\right) = \psi'(1) = 1_F = \psi\left(\frac{b}{1}\right)\psi\left(\frac{1}{b}\right)
\]
in \( F \), so
\[
\psi\left(\frac{1}{b}\right) = \psi\left(\frac{1}{b}\right),
\]
hence, for all \( a \) and \( b \) in \( R \) with \( b \neq 0 \),
\[
\psi\left(\frac{a}{b}\right) = \psi'\left(\frac{a}{1}\right)\psi'\left(\frac{1}{b}\right) = \psi\left(\frac{a}{1}\right)\psi\left(\frac{1}{b}\right) = \psi\left(\frac{a}{b}\right).
\]
This proves the second claim. Finally, we show:

**Claim 3.** The map \( i : R \to K \) satisfies the uniqueness property relative to the diagram (†).

Suppose that \( j : R \to K' \) also satisfies (†). Then there exist unique ring monomorphisms \( \psi : K \to K' \) and \( \psi' : K' \to K \) such that we have commutative diagrams
\[
\begin{array}{ccc}
R & \xrightarrow{i} & K \\
\downarrow{j} & & \downarrow{\psi} \\
K' & \xleftarrow{i} & K
\end{array}
\quad \text{and} \quad
\begin{array}{ccc}
R & \xrightarrow{j} & K' \\
\downarrow{i} & & \downarrow{\psi'} \\
K & \xleftarrow{i} & K
\end{array}
\]
By (†), we have \( \psi'\psi = 1_K \) and \( \psi' = 1_K' \) as
\[
\begin{array}{ccc}
R & \xrightarrow{i} & K \\
\downarrow{i} & & \downarrow{1_K} \\
K & \xleftarrow{i} & K
\end{array}
\quad \text{and} \quad
\begin{array}{ccc}
R & \xrightarrow{j} & K' \\
\downarrow{j} & & \downarrow{1_{K'}} \\
K' & \xleftarrow{i} & K
\end{array}
\]
It follows that \( \psi \) and \( \psi' \) are inverse isomorphisms. \( \square \)

**Notation 26.15.** If \( R \) is a domain, we shall always view \( R \subset qf(R) \), i.e., we shall identify \( R \) with its image \( \{\frac{r}{1} \mid r \in R\} \) in \( qf(R) \), e.g., \( \mathbb{Z} \subset \mathbb{Q} \).

**Remarks 26.16.** In general, it is very difficult to define a map between objects. We may have an idea where generators of the domain object must go, but to show the relations of the domain are preserved usually presents problems, i.e., to show the putative map is well-defined. Universal properties are methods to address this in certain cases. One tries to construct an object with given properties together with a map, so that there is automatically a map (that preserves the desired structure) between it and any other other object that satisfies the given properties. Moreover, that map is unique. Therefore, if you have such
a universal property, you automatically obtain unique maps to other appropriate objects. The above is an example of such. Another, that you have seen in linear algebra, is given a vector space together with a given basis, any linear transformation from it to another vector space is completely determined where the basis maps.

The following is a useful observation:

**Observation 26.17.** Let $F$ be a field and $R$ a nontrivial ring. If $\varphi : F \to R$ is a ring homomorphism, then $\varphi$ is monic as $\ker \varphi < F$ and $F$ is simple.

If $F$ is a field, then we can say a bit more.

**Corollary 26.18.** Let $F$ be a field of characteristic $n$.

1. If $n = 0$, then there exists a (unique) monomorphism $\mathbb{Q} \to F$.
2. If $n > 0$, then $n = p$ is a prime and there exists a (unique) monomorphism $\mathbb{Z}/p\mathbb{Z} \to F$.

If $K$ is a field, there exists a unique smallest subfield $\Delta_K$ of $K$, viz., the intersection of all the subfields of $K$. This field is called the prime subfield of $K$. By the corollary, we know that

$$\Delta_K \cong \begin{cases} \mathbb{Q} & \text{if } \text{char}(K) = 0. \\ \mathbb{Z}/p\mathbb{Z} & \text{if } \text{char}(K) = p. \end{cases}$$

For example, let $K$ be a finite field with $N$ elements. Then there exists a unique prime $p$ such that $\text{char} K = p$ and a unique monomorphism $\mathbb{Z}/p\mathbb{Z} \to K$. We also have $K$ is a finite dimensional vector space over $\mathbb{Z}/p\mathbb{Z}$, say of dimension $n$. So $N = p^n$. We know that $K^\times$ is a group of order $N - 1$, so $x^{N-1} = 1$ for all $x$ in $K^\times$ and $x^N = x$ for all $x \in K$.

This shows that the analogue of Fermat’s Little Theorem holds for all finite fields. We shall show later that $K^\times$ is, in fact, a cyclic group as well as showing that for every prime $p$ and positive integer $m$ there exists a field of cardinality $p^m$, and that this field is unique up to an isomorphism of fields. We also point out that there are many infinite fields of characteristic $p > 0$. Indeed, the ring $(\mathbb{Z}/p\mathbb{Z})[t]$ is an infinite domain (why?), hence so is its quotient field.

Next we turn to a generalization of the Chinese Remainder Theorem.

If $R_i, i \in I$, are rings, then the cartesian product of the $R_i$, $R = \times_I R_i$ becomes a ring via componentwise operations. We also write this as $\prod R_i$. It is commutative if all the $R_i$ are commutative. It is easy to check that $R^\times = \prod_I (R_i)^\times$. [As usual, we are sloppy about writing elements in $\prod_I R_i$.]
Theorem 26.19. (Chinese Remainder Theorem) Suppose that $R$ is a commutative ring with $\mathfrak{A}_1, \ldots, \mathfrak{A}_n$ ideals in $R$ that are comaximal, i.e., $\mathfrak{A}_i + \mathfrak{A}_j = R$ for all $i, j = 1, \ldots, n$ with $i \neq j$. Let

$$\varphi : R \to \prod_{i=1}^{n} R/\mathfrak{A}_i$$

be given by $r \mapsto (r + \mathfrak{A}_1, \ldots, r + \mathfrak{A}_n)$.

Then $\varphi$ is a ring epimorphism with

$$\ker \varphi = \mathfrak{A}_1 \cap \cdots \cap \mathfrak{A}_n = \mathfrak{A}_1 \cdots \mathfrak{A}_n,$$

hence induces isomorphisms

$$\varphi : R/(\mathfrak{A}_1 \cdots \mathfrak{A}_n) \to R/\mathfrak{A}_1 \times \cdots \times R/\mathfrak{A}_n$$

$$\varphi : (R/(\mathfrak{A}_1 \cdots \mathfrak{A}_n))^\times \to (R/\mathfrak{A}_1)^\times \times \cdots \times (R/\mathfrak{A}_n)^\times.$$

Proof. Clearly we have $\varphi$ is a ring homomorphism with $\ker \varphi = \mathfrak{A}_1 \cap \cdots \cap \mathfrak{A}_n$. By this and the First Isomorphism Theorem, we need only show that $\mathfrak{A}_1 \cap \cdots \cap \mathfrak{A}_n = \mathfrak{A}_1 \cdots \mathfrak{A}_n$ and $\varphi$ is onto, i.e., if $x_1, \ldots, x_n$ lie in $R$, there exists an $x \in R$ satisfying $x \equiv x_i \mod \mathfrak{A}_i$ for all $i$. We prove this by induction on $n$.

$n = 2$: We have $R = \mathfrak{A}_1 + \mathfrak{A}_2$, so $1 = a_1 + a_2$ for some $a_i \in \mathfrak{A}_i$, $i = 1, 2$. Therefore, $a_1 \equiv 1 \mod \mathfrak{A}_2$ and $a_2 \equiv 1 \mod \mathfrak{A}_2$; hence

$x = x_1a_2 + x_2a_1 \equiv x_i \mod \mathfrak{A}_i$ for $i = 1, 2$. As $\mathfrak{A}_1\mathfrak{A}_2 \subset \mathfrak{A}_1 \cap \mathfrak{A}_2$, we need only show that $\mathfrak{A}_1 \cap \mathfrak{A}_2 \subset \mathfrak{A}_1\mathfrak{A}_2$ to complete the $n = 2$ case. If $b \in \mathfrak{A}_1 \cap \mathfrak{A}_2$, then $b = b \cdot 1 = ba_1 + ba_2 \in \mathfrak{A}_1\mathfrak{A}_2$ since $\mathfrak{A}_1\mathfrak{A}_2 \subset \mathfrak{A}_1\mathfrak{A}_2$.

[Warning: If $\mathfrak{B}_1, \mathfrak{B}_2$ are ideals in a ring, usually $\mathfrak{B}_1\mathfrak{B}_2 < \mathfrak{B}_1 \cap \mathfrak{B}_2$.]

$n > 2$: By hypothesis, we have $R = \mathfrak{A}_1 + \mathfrak{A}_i$ for $i > 1$, so $1 = a_i + b_i$ for some $a_i \in \mathfrak{A}_1$ and $b_i \in \mathfrak{A}_i$ for each $i > 1$. Consequently, we have $1 = \prod_{i=2}^{n}(a_i+b_i)$ lies in $\mathfrak{A}_1+(\mathfrak{A}_2\cdots\mathfrak{A}_n)$ (why?), so $R = \mathfrak{A}_1+(\mathfrak{A}_2\cdots\mathfrak{A}_n)$. By the $n = 2$ case and induction, we conclude that $\mathfrak{A}_1 \cap \cdots \cap \mathfrak{A}_n = \mathfrak{A}_1 \cap (\mathfrak{A}_2\cdots\mathfrak{A}_n) = \mathfrak{A}_1 \cdots \mathfrak{A}_n$.

By induction, there exists a $y$ in $R$ satisfying

$$y \equiv x_i \mod \mathfrak{A}_i$$

for $i = 2, \ldots, n$; and by the $n = 2$ case, there exists an $x$ in $R$ satisfying

$$x \equiv x_1 \mod \mathfrak{A}_1$$

$$x \equiv y \mod \mathfrak{A}_2 \cdots \mathfrak{A}_n.$$

As $\mathfrak{A}_2\cdots\mathfrak{A}_n \subset \mathfrak{A}_i$ for $i = 2, \ldots, n$, we have $x \equiv x_i \mod \mathfrak{A}_i$ for $i = 1, \ldots, n$, proving that $\varphi$ is surjective.

We leave the isomorphism of multiplicative groups as an exercise. \qed

Exercises 26.20.
1. Let \( \varphi : R \to S \) be a ring epimorphism (of rings). Do an analogous analysis before The First Isomorphism Theorem for Groups 12.1 for \( \varphi : R \to S \) to see how ideals and factor rings arise naturally.

2. Let \( R = (\mathbb{Z}/2\mathbb{Z})[t] \), \( f = f(t) = t^2 + t + 1 \), and \( g = t^2 + 1 \). Show all of the following:
   (i) \( R/(f) \) is a field with four elements.
   (ii) \( R/(g) \) is not a domain and has four elements.
   (iii) Neither \( R/(f) \) nor \( R/(g) \) is isomorphic to the ring \( \mathbb{Z}/4\mathbb{Z} \).

3. Let \( R \) be a commutative ring. Suppose for every element \( x \) in \( R \) there exists an integer \( n = n(x) > 1 \) such that \( x^n = x \). Show that every prime ideal in \( R \) is maximal.

4. (Second Isomorphism Theorem) Let \( R \subset S \) be a subring, \( \mathfrak{A} \subset S \) an ideal. Then \( \mathfrak{A} \cap R \) is an ideal in \( R \), \( R + \mathfrak{A} \) is a subring of \( S \), and \( (R + \mathfrak{A})/\mathfrak{A} \cong R/(R \cap \mathfrak{A}) \).

5. Show that the inclusion map \( i : \mathbb{Z} \subset \mathbb{Q} \) satisfies the property that if \( \psi_1 \circ i = \psi_2 \circ i \), for \( \psi_1, \psi_2 : \mathbb{Q} \to R \) ring homomorphisms, then \( \psi_1 = \psi_2 \). This shows that a surjective ring homomorphism is not equivalent to this property. (Cf. Remark 9.5.)

6. Let \( R \) be a commutative ring and \( A = R[[t]] \), the power series over \( R \). Show if \( \mathfrak{A} \) is an ideal in \( R \), then \( \mathfrak{A} \) is a prime ideal in \( R \) if and only \( A\mathfrak{A} \) is a prime ideal in \( A \). Moreover, if \( \mathfrak{m} \) is a maximal ideal in \( R \), then \( \mathfrak{A} \mathfrak{m} + At \) is a maximal ideal in \( A \) and the unique maximal ideal in \( A \) containing \( \mathfrak{m} \). In particular, if \( R \) has precisely one maximal ideal, then so does \( A \). [A commutative ring having precisely one maximal ideal is called a local ring.]

7. Let \( R \) be a commutative ring of characteristic \( p > 0 \), \( p \) a prime. Prove that the map \( R \to R \) by \( x \mapsto x^p \) is a ring homomorphism. It is called the Frobenius homomorphism. In particular, the Children’s Binomial Theorem holds, i.e., \( (x + y)^p = x^p + y^p \) in \( R \) for all \( x \) and \( y \) in \( R \).

8. Let \( R \) be a finite domain. Show that \( \text{char}(R) = p \) is a prime and \( R \) is a field. Moreover, \( R \) is a vector space over \( \mathbb{Z}/p\mathbb{Z} \).

9. Show that if \( R \) is a domain, so is the polynomial ring \( R[t] \). In particular, show that there exists fields properly containing the complex numbers. Does the field that you constructed have the property that every non-constant polynomial over it has a root? Prove or disprove this.

10. Let \( \varphi : R \to S \) be a ring isomorphism. Show \( \varphi \) induces a group isomorphism \( R^\times \to S^\times \) by restriction. We also write this map as \( \varphi \).
11. Prove the isomorphism statement about the multiplicative groups in the Chinese Remainder Theorem 26.19

27. Zorn’s Lemma

One would like to have a method of proof generalizing induction. There are many equivalent ways of doing this: Well-Ordering Principle, Transfinite Induction, Zorn’s Lemma. These are also equivalent to the Axiom of Choice, Tychonoff’s Theorem, the theorem that all vector spaces have bases. In algebra, the most useful form is Zorn’s Lemma, which generalizes the Well-Ordering Principle. These are independent of the other axioms in Zermelo-Frankel Set Theory, and we accept them as an axiom. [In Appendix A, we show that the Axiom of Choice, Zorn’s Lemma, and the Well-Ordering Principle are all equivalent.]

Definition 27.1. Let $S$ be a set with a relation $R$ on $S$. We call $S$ a partially ordered set or poset (under $R$) if the following hold: For all $a, b, c$ in $S$

1. $aRa$. [Reflexitivity]
2. If $aRb$ and $bRa$ then $a = b$.  
3. If $aRb$ and $bRc$ then $aRc$. [Transitivity]

We also say $(S, R)$ is a poset, if the relation is not obvious. So the second condition replaces symmetry in the definition of an equivalence relation. If $S$ is a poset (under $R$), $a, b$ elements in $S$, then $a$ and $b$ may be incomparable, i.e., neither $aRb$ nor $bRa$. We say that the poset $S$ is a chain or totally ordered if for all $a$ and $b$ in $S$, we have either $aRb$ or $bRa$, i.e., all elements are comparable.

Examples 27.2. 1. Recall if $T$ is a set, then the power set of $T$ is defined by $\mathcal{P}(T) := \{A \mid A \subset T\}$. It becomes a poset under $\subset$ (or $\supset$). It is almost never a chain. [When is it?]  
2. $(\mathbb{Z}, \leq)$ and $(\mathbb{R}, \geq)$ are chains.  
3. As $\{1, \sqrt{-1}\}$ is a basis for $\mathbb{C}$ as a vector space over $\mathbb{R}$, for each $\alpha \in \mathbb{C}$, there exist unique $x, y \in \mathbb{R}$ such that $\alpha = x + y\sqrt{-1}$. Define the lexicographic order on $\mathbb{C}$ by $\alpha = x + y\sqrt{-1} \leq_L \beta = w + z\sqrt{-1}$ if $\alpha = \beta$ and if not this, then $x < w$ and if not these, then $(x = w$ and $y < z$) in $\mathbb{R}$. Then $(\mathbb{C}, \leq_L)$ is a chain.  

Note: If $\alpha \leq_L \beta$ and $\gamma \leq_L \delta$, then $\alpha + \gamma \leq_L \beta + \delta$, but $\alpha \gamma >_L \beta \delta$ is possible (where $>_L$ is the obvious relation), e.g., $0 \leq_L \sqrt{-1}$ but $\sqrt{-1}\sqrt{-1} <_L 0$. 

[One can generalize the lexicographic order on $\mathbb{C}$ to any finite dimensional vector space $V$ over $\mathbb{R}$ with ordered basis $\{v_1, \ldots, v_n\}$. How?]

**Definition 27.3.** Let $(S, R)$ be a poset, $T$ a subset of $S$. An element $a$ in $S$ is called an *upper bound* of $T$ if $xRa$ for all $x \in T$, and $s$ is called a *maximal element* of $S$ if $sRy$ with $y \in S$ implies that $s = y$. We say that the poset $S$ is *inductive* if every chain in $S$ has an upper bound in $S$.

**Examples 27.4.** 1. $1$ is a maximal element in $([0, 1], \leq)$.
2. $0$ is a maximal element in $([0, 1], \geq)$. (Of course, it is usually called a *minimal element*.)
3. $(\mathbb{Z}, \leq)$ is not inductive.
4. $(\mathbb{Z}^+, \geq)$ is inductive. (Note here ‘upper bound’ is usually called ‘lower bound’ for obvious reasons.)

We can now state the axiom that we shall use as an extension of finite induction.

**Lemma 27.5.** *(Zorn’s Lemma)* Let $S$ be a non-empty inductive poset. Then $S$ contains a maximal element.

This is indeed an extension of finite induction, for if $S$ is a subset of $\mathbb{Z}^+$, then $(S, \geq)$ has a maximal element, which is exactly the Well-Ordering Principle. As mentioned, we cannot prove this lemma; we accept it as an axiom.

We give some applications of Zorn’s Lemma. The first is the result that vector spaces always have bases.

**Proposition 27.6.** Let $V$ be a nonzero vector space over a field $F$ and $S$ a linearly independent subset of $V$. Then $S$ extends to a basis of $V$, i.e., is part of a basis for $V$.

**Proof.** Let

$$S = \{ T \mid S \subset T \subset V \text{ with } T \text{ linearly independent} \},$$

a nonempty poset under $\subset$. Let $C$ be a chain in $S$.

**Claim.** $A = \bigcup_C T$ is an upper bound for $C$ in $S$:

Suppose that $A$ is linearly dependent. By the definition of linear dependence, there exists a finite linearly dependent subset $\{v_1, \ldots, v_n\}$ of $A$. For each $i$, $1 \leq i \leq n$, there exists a $T_i \in C$ with $v_i \in T_i$. Since $C$ is a chain and $n$ finite, there exists a $T$ in $C$ containing $T_i$ for $i = 1, \ldots, n$, hence containing $\{v_1, \ldots, v_n\}$, contradicting the fact that $T$ is linearly independent. This shows that $A$ is indeed an upper bound of $C$ in $S$. 

As the hypotheses of Zorn’s Lemma have been fulfilled, there exists a maximal element \( T \) in \( S \). We show that \( T \) is a basis for \( V \). If the span, \( \langle T \rangle \), of \( T \) is not \( V \), then there exists \( w \in V \setminus \langle T \rangle \). Let \( T' = T \cup \{ w \} \). As \( T' \) must be linearly dependent by maximality of \( T \), there exist \( \beta \) and \( \beta_v \) in \( F \) for \( v \in T \), with almost all \( \beta_v = 0 \), i.e., \( \beta_v = 0 \) except for finitely many \( v \) in \( T \), but not all \( \beta, \beta_v \) zero that satisfies \( \sum_{T} \beta_v v + \beta w = 0 \). As \( T \) is linearly independent, \( \beta \neq 0 \), so \( w = -\sum_{T} \beta^{-1} \beta_v v \) lies in \( \langle T \rangle \), a contradiction. Consequently, \( T \) spans \( V \), so is a basis for \( V \). \( \square \)

**Proposition 27.7.** Let \( V \) be a nonzero vector space over a field \( F \) and \( S \) a spanning set for \( V \). Then a subset of \( S \) is a basis for \( V \).

We leave a proof of this as an exercise. You should try the obvious proof. It does not work. Can you figure out what this means? Rather you will have to modify the proof above.

The two proposition imply, just as in the finite dimensional case, the following:

**Corollary 27.8.** Let \( V \) be a nonzero vector space over a field \( F \) and \( B \) a subset of \( V \). Then the following are equivalent:

1. \( B \) is a basis for \( V \).
2. \( B \) is a maximal linearly independent set in \( V \) (relative to \( \subset \)).
3. \( B \) is a minimal spanning set for \( V \) (relative to \( \supset \)).

As a further application of Zorn’s Lemma, we next give a proof of the crucial result Theorem 15.13 needed in Section 15. The proof is a typical Zorn’s Lemma argument used to extend maps.

**Proposition 27.9.** Let \( G \) be a divisible additive group and \( \theta : A \to B \) a group monomorphism of additive groups. If \( \varphi : A \to G \) is a group homomorphism, then there exists a group homomorphism \( \psi : B \to G \) satisfying \( \varphi = \psi \circ \theta \).

**Proof.** We may replace \( A \) by \( \theta(A) \), i.e., assume that \( A \subset B \). Let \( S = \{(C, \psi_C) \mid A \subset C \subset B, \psi_C : C \to G \ \text{a group homomorphism with} \ \varphi_C|_A = \varphi\} \).

Partially order \( S \) by

\[
(C, \psi_C) \leq (C', \psi'_{C'}) \iff C \subset C' \ \text{and} \ \psi_C|_C = \psi_{C'}.
\]

The set \( S \) is not empty, since \( (A, \varphi) \in S \). Let \( \mathcal{C} \) be a chain in \( S \) and set \( C = \bigcup_{(C, \psi_C) \in \mathcal{C}} C_{a} \). Define \( \psi_C : C \to G \) by \( x \mapsto \psi_{C_a}(x) \) if \( x \in C_a \). Since \( \mathcal{C} \) is a chain, \( \psi_C \) is well-defined. Therefore, \( (C, \psi_C) \in S \) is an upper bound for \( \mathcal{C} \). By Zorn’s Lemma, there exist a maximal element \( (C, \psi_C) \in S \).
Claim. \( C = B \):
Suppose the claim does not hold. Then there exists an \( x \in B \setminus C \).

Case 1. \( (x) \cap C = 0 \):
In this case \( C + \langle x \rangle = C \oplus \langle x \rangle \). Define \( \widetilde{\psi} = \psi \oplus 0 \), i.e., \( (c, rx) \mapsto \psi_C(c) \).
Then \( (C, \psi_C) < (C \oplus \langle x \rangle, \widetilde{\psi}) \) in \( S \), contradicting the maximality of \( (C, \psi) \).

Case 2. \( (x) \cap C \neq 0 \):
Choose \( n \in \mathbb{Z}^+ \) minimal such that \( nx \in C \). Since \( G \) is divisible, there exists \( y \in G \) satisfying \( ny = \psi_C(nx) \). Define the map \( \widetilde{\psi} : C + \langle x \rangle \to G \) defined by \( c + rx \mapsto \psi_C(c) + ry \). Then \( \widetilde{\psi} \) is a group homomorphism extending \( \psi \) to \( C + \langle x \rangle \). As \( (C, \psi_C) < (C + \langle x \rangle, \widetilde{\psi}) \) in \( S \), this contradicts the maximality of \( (C, \psi) \). \( \square \)

We next apply Zorn’s Lemma to Ring Theory.

**Proposition 27.10.** Let \( R \) be a nontrivial ring and \( \mathfrak{A} \subset R \) an ideal.
Then there exists a maximal ideal \( \mathfrak{m} \) in \( R \) such that \( \mathfrak{A} \subset \mathfrak{m} \). In particular, any nontrivial ring contains maximal ideals.

**Proof.** Let
\[ S = \{ \mathfrak{B} | \mathfrak{A} \subset \mathfrak{B} < R \text{ is an ideal} \}. \]
The set \( S \) is a poset under \( \subset \), and \( S \) is non-empty as \( \mathfrak{A} \in S \). Let \( C \) be a chain in \( S \), so
for all \( \mathfrak{B} \) and \( \mathfrak{B}'' \) in \( C \) either \( \mathfrak{B}' \subset \mathfrak{B}'' \) or \( \mathfrak{B}'' \subset \mathfrak{B}' \).

**Claim.** \( \bigcup_{\mathfrak{B} \in C} \mathfrak{B} < R \) is an ideal and hence an upper bound for \( C \) in \( S \):
Let \( r \in R \) and \( x, y \in \bigcup_{\mathfrak{B} \in C} \mathfrak{B} \). Then there exist \( \mathfrak{B}' \) and \( \mathfrak{B}'' \) in \( C \) with \( x \in \mathfrak{B}' \in C \) and \( y \in \mathfrak{B}'' \in C \). As \( C \) is a chain, we may assume that \( \mathfrak{B}' \subset \mathfrak{B}'' \), so \( rx, xr, x \pm y \in \mathfrak{B}' \subset \mathfrak{B}'' \subset \bigcup_{\mathfrak{B} \in C} \mathfrak{B} \), showing \( \bigcup_{\mathfrak{B} \in C} \mathfrak{B} \) is an ideal. If \( \bigcup_{\mathfrak{B} \in C} \mathfrak{B} \) is the unit ideal, there exists a \( \mathfrak{B}' \) in \( C \) such that \( 1 \in \mathfrak{B}' \). This means that \( \mathfrak{B}' \) is the unit ideal, so not in \( S \), a contradiction. By Zorn’s Lemma, there exists an ideal \( \mathfrak{m} \in S \) that is a maximal element. So \( \mathfrak{m} \) satisfies \( \mathfrak{A} \subset \mathfrak{m} < R \) is an ideal, maximal with respect to this property. Thus \( \mathfrak{m} \) is a maximal ideal of \( R \) containing \( \mathfrak{A} \). \( \square \)

**Remark 27.11.** The reason the above worked was the claim, i.e.,

(i) The union of ideals in a chain is an ideal under the partial order \( \subset \).

[The same is true for chains of subgroups, subspaces of a vector space, and other algebraic objects.]

**Note.** This does not work for arbitrary unions.
(ii) There exists an element lying in none of the elements in any chain (in this case the element 1).
In fact, if $R$ is a non-trivial rng (even commutative), it may not have maximal ideals.

We next want to give further application of Zorn’s Lemma to commutative ring theory that generalizes the existence of maximal ideals and which is quite useful. We begin with an important definition.

**Definition 27.12.** Let $R$ be a nontrivial commutative ring, $S$ a nonempty subset of $R$. We call $S$ a **multiplicative set** if

1. $1 \in S$.
2. If $s_1$ and $s_2$ are elements of $S$, then $s_1 s_2$ lies in $S$, i.e., $S$ is closed under multiplication.

If $T$ is a non-empty subset of $R$, we say that $T$ **excludes** $S$ if $S \cap T = \emptyset$.

**Examples 27.13.** Let $R$ be a nontrivial commutative ring, $a \in R$, and $\mathfrak{A} < R$ an ideal.

1. $\mathfrak{A}$ is a prime ideal if and only if $R \setminus \mathfrak{A}$ is a multiplicative set.
2. $\{a^n \mid n \geq 0\}$ is a multiplicative set.
3. Let $S$ be a multiplicative set with $0 \notin S$. Then $S$ contains no nilpotent elements, i.e., an element $x \in R$ that satisfies $x^n = 0$ for some $n \in \mathbb{Z}^+$.
4. $R^\times$ is a multiplicative set.
5. An element $x \in R$ is called a **zero divisor** of $R$ if there exists a nonzero element $y \in R$ satisfying $xy = 0$. Let $zd(R) := \{x \in R \mid x$ is a zero divisor}\}. Then $R \setminus zd(R)$ is a multiplicative set.

Using the concept of multiplicative sets, we can now prove the generalization of the existence of maximal ideals that we seek. The proof will be just as the one for the existence of maximal ideals but using the ideal analogue for prime ideals of Euclid’s Lemma given in Lemma 25.20.

**Theorem 27.14.** (Krull) Let $R$ be a commutative ring and $S$ a multiplicative set. Suppose that $\mathfrak{A} < R$ is an ideal in $R$ excluding $S$. Then there exists an ideal $\mathfrak{B} \subset R$ containing $\mathfrak{A}$ excluding $S$ and maximal with respect to these two properties. Moreover, any such $\mathfrak{B}$ is a prime ideal.
**Note:** If $R$ is a non-trivial commutative ring, applying the above to the multiplicative set $R^\times$ shows the existence of maximal ideals in a commutative ring.

**Proof.** Let 

$$I = \{ \mathfrak{B} \mid \mathfrak{A} \subset \mathfrak{B} < R \text{ an ideal excluding } S \}$$

We know that $I$ is not empty as $\mathfrak{A} \in I$. Partially order $I$ by $\subset$. Let $\mathcal{C}$ be a chain in $I$.

**Claim.** $\bigcup \mathcal{C}$ is an ideal in $R$ excluding $S$ (hence not the unit ideal).

In particular, $\bigcup \mathcal{C}$ is an upper bound for $\mathcal{C}$ in $I$:

$$\bigcup \mathcal{C} \text{ is an ideal exactly as before. If } b \in S \cap \bigcup \mathcal{C}, \text{ then there exists an ideal } \mathfrak{B} \in \mathcal{C} \text{ such that } b \in S \cap \mathfrak{B}, \text{ a contradiction. This establishes the claim.}$$

By Zorn’s Lemma, there exists an ideal $\mathfrak{B}$ in $I$ satisfying

1. $\mathfrak{B} \cap S = \emptyset$.
2. $\mathfrak{B} \subset \mathfrak{B}$. (iii) If $\mathfrak{B} < \mathfrak{B} \subset R$ is an ideal, then $\mathfrak{B} \cap S \neq \emptyset$.

To finish, we must show that $\mathfrak{B}$ is a prime ideal. Suppose this is not the case. Then by Lemma 25.20, there exist ideals $\mathfrak{B} < \mathfrak{B}, i = 1, 2$, satisfying $\mathfrak{B} \mathfrak{B} \subset \mathfrak{B}$. By (iii), there exist $s_i \in \mathfrak{B} \cap S, i = 1, 2$. As $S$ is a multiplicative set, $s_1s_2 \in \mathfrak{B} \cap S$, a contradiction. $\Box$

Note the last part of the proof is structurally similar to factoring integers into primes. This is an argument that is often used.

**Example 27.15.** Let $C([0, 1])$ denote the ring of continuous real-valued functions on $[0, 1]$. Exercise 25.21(13) said that the maximal ideals in $C([0, 1])$ are in one to one correspondence with $[0, 1]$, the correspondence given by $a \rightarrow m_a = \{ f \in C([0, 1]) \mid f(a) = 0 \}$. [This is true with $[0, 1]$ replaced by any finite closed interval (or any compact subset) of the real line.] Suppose instead, we look at the ring $C(0, 1)$, the continuous real-valued functions on the open interval $(0, 1)$ (or on the half closed interval $C(0, 1)$). Let $S$ denote the set of all nonzero polynomial functions in $C(0, 1)$. (Cf. Remark 33.10(3) for a precise definition of a polynomial function.) Then $S$ is a multiplicative set excluding the zero ideal. Hence by the theorem, there exists a prime ideal $\mathfrak{p}$ excluding $S$. Consequently, $\mathfrak{p}$ does not contain any linear polynomial function $x - a, a \in (0, 1)$. In particular, $\mathfrak{p} \not\subset m_a$ for any $a \in (0, 1)$, so there exist maximal ideals not of the form $m_a$.

The theorem allows us to begin to see the importance of prime ideals in understanding the structure of a commutative ring.
Definition 27.16. Let $R$ be a non-trivial commutative ring. The nilradical of $R$ is the set
\[
\text{nil}(R) := \{a \mid a \text{ is a nilpotent element in } R\}.
\]

Corollary 27.17. Let $R$ be a non-trivial commutative ring. Then \(\text{nil}(R)\) is an ideal in $R$ and
\[
\text{nil}(R) = \bigcap_{p \text{ a prime ideal}} p.
\]

Proof. We leave it as an exercise to prove \(\text{nil}(R)\) is a ideal.

\[
\text{nil}(R) \subseteq \bigcap_{p \text{ a prime ideal}} p: \text{ If } a \in \text{nil}(R), \text{ there exists an } n \in \mathbb{Z}^+ \text{ such that } a^n = 0 \text{ and } 0 \in \mathfrak{a} \text{ for every ideal, in particular for prime ideals. Hence } a \in p \text{ for all prime ideals } p \text{ in } R.
\]

\[
\bigcap_{p \text{ a prime ideal}} p \subseteq \text{nil}(R): \text{ Suppose this is not true, then there exists an element } a \in \bigcap_{p \text{ a prime ideal}} p \setminus \text{nil}(R).
\]

In particular, the multiplicative set $S = \{a^n \mid n \geq 0\}$ does not contain 0 hence the trivial ideal (0) excludes $S$. By the theorem of Krull, there exists a prime ideal $p$ excluding $S$. Thus $a \notin p$, a contradiction.

Therefore, \(\text{nil}(R) = \bigcap_{p \text{ a prime ideal}} p\) as needed. \(\square\)

Note. If $R$ is a domain, then \(\text{nil}(R) = 0\). The converse is false, e.g., $\mathbb{Z} \times \mathbb{Z}$ has trivial nilradical, but is not a domain. A commutative ring whose nilradical is trivial is called a reduced ring.

Remark 27.18. As the intersection of a chain of prime ideals in a nonzero commutative ring $R$ is clearly a prime ideal, there exist minimal elements in the set of all prime ideals by Zorn’s Lemma. We call such a prime a minimal prime ideal of $R$. As every prime ideal in $R$ contains a minimal prime, we have
\[
\text{nil}(R) = \bigcap_{p \text{ a minimal prime ideal}} p.
\]

Exercises 27.19.
1. Let $V$ be a finite dimensional vector space over $\mathbb{R}$ with ordered basis $\{v_1, \ldots, v_n\}$. Define a lexicographic order of $V$ relative to this ordered basis.

2. Prove Proposition 27.7.

3. Rngs may not have maximal ideals. This problem constructs one. Let $p$ be a prime number. Let $\mathbb{Z}/p\mathbb{Z}^\infty$ be the (additive) subgroup of $\mathbb{Q}/\mathbb{Z}$ consisting of all elements having order some power of $p$, i.e., $x = \alpha + \mathbb{Z}$, $\alpha \in \mathbb{Q}$, lies in $\mathbb{Z}/p\mathbb{Z}^\infty$ if and only if $p^r x = 0$ in $\mathbb{Q}/\mathbb{Z}$, i.e., $p^r \alpha$ is an integer, for some positive integer $r$.
   
   (i) Show that the set $\{\frac{1}{p^r} + \mathbb{Z} | r \text{ a non-negative integer} \}$ generates $\mathbb{Z}/p\mathbb{Z}^\infty$.

   (ii) Show the subgroup $\langle \frac{1}{p^r} + \mathbb{Z} \rangle$ of $\mathbb{Z}/p\mathbb{Z}^\infty$ is isomorphic to $\mathbb{Z}/p^r \mathbb{Z}$.

   (iii) Show that any subgroup of $\mathbb{Z}/p\mathbb{Z}^\infty$ is $\langle \frac{1}{p^r} + \mathbb{Z} \rangle$ for some non-negative integer $r$ and the subgroups of $\mathbb{Z}/p\mathbb{Z}^\infty$ form a chain under set inclusion.

   (iv) Make $\mathbb{Z}/p\mathbb{Z}^\infty$ into a rng by defining $x \cdot y = 0$ for all $x, y \in \mathbb{Z}/p\mathbb{Z}^\infty$. Show that $\mathbb{Z}/p\mathbb{Z}^\infty$ has no maximal ideals.

4. Let $R$ be a commutative ring. Prove the following:
   
   (i) Let $S$ be the set of non finitely generated ideals in $R$ and suppose that it is not empty. Let $\mathfrak{A}$ be a maximal element in $S$. Then $\mathfrak{A}$ is a prime ideal.

   (ii) Let $\mathfrak{A}$ be a non finitely generated ideal in $R$. Suppose an ideal $\mathfrak{B}$ in $R$ has the property that it contains $\mathfrak{A}$, is not finitely generated, and is maximal with respect to this property. Show that there exists a prime ideal in $R$ containing $\mathfrak{A}$ that is not finitely generated.

5. Let $R$ be a commutative ring and $A = R[[t]]$, the power series over $R$. Let $\varphi : A \to R$ be the ring epimorphism induced by $t \mapsto 0$ (evaluation at 0). Show the following:
   
   (i) If $\mathfrak{A}$ is an ideal in $R$, then $\varphi^{-1}(\mathfrak{A}) = \mathfrak{A}A + tA$.

   (ii) $\varphi$ induces a bijection between the set of maximal ideals in $A$ and the set of maximal ideals in $R$.

6. As in the previous problem, let $R$ be a commutative ring, $A = R[[t]]$, and $\varphi : A \to R$ be the ring epimorphism induced by $t \mapsto 0$. Show the following:
   
   (i) Let $\mathfrak{P}$ be a prime ideal in $A$. Suppose that $\varphi(\mathfrak{P}) = (a_1, \ldots, a_n)$. Then $\mathfrak{P} = Aa_1 + \cdots + Aa_n + At$ if $t \in \mathfrak{P}$ and $\mathfrak{P} = (f_1, \ldots, f_n)$ for $f_i \in A$ with $f_i = a_i +\text{ higher terms in } t \text{ if } t \notin \mathfrak{P}$.
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(ii) Every prime ideal in \( A \) can be generated by finitely many elements if and only if every prime ideal in \( R \) can be generated by finitely many elements.

7. Let \( R \) be a commutative ring. Prove
   
   (i) If \( x \) is nilpotent in \( R \), then \( 1 + x \) is a unit in \( R \).
   
   (ii) The nilradical of \( R \) is an ideal.
   
   (iii) Compute the nilradical of the rings: \( \mathbb{Z}/12\mathbb{Z}, \mathbb{Z}/n\mathbb{Z}, n > 1, \) and \( \mathbb{Z} \).

8. Let \( R \) be a commutative ring. The Jacobson radical of \( R \) is defined to be \( \text{rad}(R) := \bigcap_{\text{Max}(R)} m \), the intersection of all maximal ideals in \( R \). Show that \( x \) lies in \( \text{rad}(R) \) if and only if \( 1 - yx \) is a unit in \( R \) for all \( y \) in \( R \).

9. Let \( R \) be a commutative ring and \( \mathfrak{A} \triangleleft R \) an ideal. Define the radical of \( \mathfrak{A} \) to be the set
   
   \[ \sqrt{\mathfrak{A}} := \{ x \in R \mid x^n \in \mathfrak{A} \text{ for some } n \in \mathbb{Z}^+ \}. \]

   Show the following:
   
   (i) \( \mathfrak{A} \) is an ideal and
   
   \[ \sqrt{\mathfrak{A}} = \bigcap_{\mathfrak{p} \triangleleft \mathfrak{A} \triangleleft R, \mathfrak{p} \text{ a prime ideal}} \mathfrak{p}. \]

   (ii) Let \(- : R \to R/\mathfrak{A}\) be the canonical ring epimorphism, then \( \text{nil}(R) = \sqrt{\mathfrak{A}}/\mathfrak{A} \).

10. Let \( R \) be a commutative ring, \( \mathfrak{A} \triangleleft R \) an ideal, and \(- : R \to R/\mathfrak{A}\) be the canonical epimorphism. We say that \( \mathfrak{A} \) is a primary ideal if \( ab \in \mathfrak{A} \) implies that \( a \in \mathfrak{A} \) or \( b^n \in \mathfrak{A} \) for some positive integer \( n \). Let \( \mathfrak{A} \triangleleft R \) be an ideal. Show both of the following:

    (i) \( \mathfrak{A} \) is a primary ideal if and only if every zero divisor of \( R/\mathfrak{A} \) is nilpotent.

    (ii) If \( \mathfrak{A} \) is primary, then its radical, \( \sqrt{\mathfrak{A}} \), (see the previous exercise) is a prime ideal.

11. Determine all primary ideals if \( R \) is a PID. (Cf. the previous exercise for the definition of a primary ideal.)

12. Let \( R \) be a commutative ring. An element \( e \) of \( R \) is called an idempotent if \( e^2 = e \). For example, if \( S \) is another commutative ring the element \((1_R, 0_S)\) is an idempotent in the ring \( R \times S \). The object of this exercise is to prove a converse. Let \( e \) be an idempotent of \( R \). Then prove

    (i) \( e' := 1 - e \) is an idempotent of \( R \).
(ii) The principal ideal $Re$ of $R$ is a ring with identity $1_{Re} = e$. [Note that $Re$ is not a subring of $R$ since $Re$ will not have the same identity as $R$ unless $e = 1$.]

(iii) $R$ is ring isomorphic to $Re \times Re'$.

28. Localization

In section 26, we showed that a field of quotients exists for any domain by defining fractions. The purpose was to create a ring containing the given ring $R$ in which every nonzero element of $R$ has a multiplicative inverse. [Note if 0 has an inverse then the ring must be trivial.] Instead, we can try to invert only some elements. This turns out to be a crucial idea, and the method is to invert elements in a multiplicative set. But to make it effective, we must weaken what we mean when we say two ‘fractions’ are equal.

Construction 28.1. We shall leave most of the details of this construction as an exercise, an exercise that you should do. Let $R$ be a commutative ring and $S$ a multiplicative set in $R$. Let

$S = \{(a, s) \mid a \in R, s \in S\} = R \times S.$

Define a relation on $S$ as follows: Write

$(a, s) \sim (a', s')$ in $S$

if there exists an $s'' \in S$ satisfying $s''(as' - a's) = 0$ in $R$.

Then $\sim$ is a equivalence relation. Let

$\frac{a}{s}$ denote the equivalence class of $(a, s)$

and set

$S^{-1}R = \{\frac{a}{s} \mid a \in R, s \in S\} = S/\sim.$

Define

$\frac{a}{s_1} + \frac{b}{s_2} := \frac{as_2 + bs_1}{s_1s_2}$

$\frac{a}{s_1} \cdot \frac{b}{s_2} := \frac{ab}{s_1s_2}$

for all $a, b \in S$ and for all $s_1, s_2 \in S$. Then

$+: S^{-1}R \times S^{-1}R \to S^{-1}R$ and $\cdot : S^{-1}R \times S^{-1}R \to S^{-1}R$

are well-defined maps making $S^{-1}R$ into a commutative ring, called the localization of $R$ at $S$ with $0_{S^{-1}R} = \frac{0}{s}$ and $1_{S^{-1}R} = \frac{s}{s}$, for any $s \in S$. 

The map

\[ \varphi : R \to S^{-1}R \] given by \( r \mapsto \frac{r}{1} \)

is a ring homomorphism.

**Examples 28.2.** Let \( R \) be a commutative ring and \( S \) a multiplicative set in \( R \).

1. If \( 0 \in S \), then \( S^{-1}R \) is the trivial ring.
2. Suppose that \( R \) is a domain. Then \( T = R \setminus \{0\} \) is a multiplicative set and \( T^{-1}R \) is the quotient field of \( R \).
3. Let \( \varphi : R \to S^{-1}R \) be the ring homomorphism given by \( r \mapsto \frac{r}{1} \). Then
   (i) \( \varphi(S) \subset (S^{-1}R)^{\times} \).
   (ii) \( \ker \varphi = \{ x \in R \mid \text{there exists an } s \in S \text{ such that } sx = 0 \} \).
   (iii) \( \varphi \) is a ring monomorphism if and only if \( S \cap \text{zd}(R) = \emptyset \).
   (iv) If \( R \) is a domain and \( S \) does not contain zero, then \( \varphi \) is a ring monomorphism. Note if this is the case that \( \frac{a}{s} = \frac{b}{s} \) in \( S^{-1}R \) if and only if \( a = b \). We then view \( \varphi \) as the inclusion map, so \( R \subset S^{-1}R \subset qf(R) \) (cf. (2)).
4. Let \( p \) be a prime ideal in \( R \). Then \( T = R \setminus p \) is a multiplicative set and the localization \( T^{-1}R \) is denoted \( R_p \). For example, if \( p \) is a prime in \( \mathbb{Z} \) and \( T = \{ n \in \mathbb{Z} \mid p \not| n \} \), then \( \mathbb{Z}_{(p)} = \{ \frac{n}{m} \mid n, m \in \mathbb{Z} \text{ relatively prime with } p \not| m \} \).
5. Let \( f \) be an element of \( R \). Then \( T = \{ f^n \mid n \geq 0 \} \) is a multiplicative set and the localization \( T^{-1}R \) is denoted \( R_f \). For example, if \( p \) is a prime in \( \mathbb{Z} \), then \( \mathbb{Z}_p = \{ \frac{n}{m} \mid n, m \in \mathbb{Z} \text{ relatively prime with } m = p^k \text{ some } k \geq 0 \} \).

The existence of the localization of a commutative ring at a multiplicative set is a basic construction in commutative algebra whose importance cannot be overestimated. Indeed, it, together with the construction of polynomial rings and quotient rings over a commutative ring are the basic constructions in commutative algebra. It is especially important when the multiplicative set is the complement of a prime ideal. For example, the ring \( \mathbb{Z}_{(p)} \) in Example 28.2(4) has only one nonzero prime ideal and all prime integers relatively prime to \( p \) become units. One can view \( \mathbb{Z}_{(p)} \subset \mathbb{Q} \). It turns out that in this ring every element can be written \( up^n \) for some unit \( u \) in \( \mathbb{Z}_{(p)} \) and some positive integer \( n \). In the more general case that \( R \) is a commutative ring and \( p \) a prime ideal in \( R \), let \( pR_p \) denote the image of \( p \) in \( R_p \). Then \( pR_p \) is the unique maximal ideal in \( R_p \). One can obtain valuable information for \( R \) from \( R_p \) via the natural map \( R \to R_p \) just as one get
valuable information for $R$ from the canonical epimorphism $R \to R/p$. Localization serves as the algebraic analogue of the geometric study of germs of functions at a point and of functions in neighborhoods of a point.

**Exercises 28.3.**

1. Prove all of the claims in Construction 28.1.

2. Let $R$ be a commutative ring and $S$ a multiplicative set in $R$. Let $\varphi : R \to S^{-1}R$ be given by $r \mapsto \frac{r}{1}$. Show that this satisfies the following universal property: If $\psi : R \to R'$ is a ring homomorphism with $R'$ commutative and $\psi(S)$ a subset of the unit group of $R'$, then there exists a unique ring homomorphism $\theta : S^{-1}R \to R'$ such that

$$
\begin{array}{ccc}
R & \xrightarrow{\varphi} & S^{-1}R \\
\downarrow{\psi} & & \downarrow{\theta} \\
R' & & \end{array}
$$

commutes.

3. Let $R$ be a commutative ring, $S$ a multiplicative set in $R$. Show if $p$ is a prime ideal in $R$ excluding $S$, then $S^{-1}p$ is a prime ideal, and every prime in $S^{-1}R$ is of this form. Use this to show if $0 \notin S$, then there exists a prime ideal in $R$ excluding $S$. (Cf. this to Theorem 27.14.)

4. Let $p$ be a prime ideal in a commutative ring $R$. Show that $R_p$ has a unique maximal ideal. What is it? A commutative ring with a unique maximal ideal is called a *local ring*.

5. Let $\mathfrak{A}$ be an ideal in a commutative ring $R$ and $S$ a multiplicative set in $R$. Define

$$
S^{-1}\mathfrak{A} := \{ \frac{a}{s} \mid a \in \mathfrak{A}, s \in S \}.
$$

Show that $S^{-1}\mathfrak{A}$ is an ideal in $S^{-1}R$ and not the unit ideal if $\mathfrak{A}$ excludes $S$.

[Note. If $\frac{x}{s}$ lies in $S^{-1}\mathfrak{A}$, this does not mean that $x \in \mathfrak{A}$, only that there exist an $a \in \mathfrak{A}$ and an $s' \in S$ such that $\frac{x}{s} = \frac{a}{s'}.$]

6. Let $R$ be a commutative ring, $S$ a multiplicative set in $R$, and $\varphi : R \to S^{-1}R$ the ring homomorphism given by $r \mapsto \frac{r}{1}$. Suppose that $\mathfrak{B}$ is an ideal in $S^{-1}R$. Show that $\mathfrak{A} = \varphi^{-1}(\mathfrak{B})$ is an ideal in $R$ and $S^{-1}\mathfrak{A} = \mathfrak{B}$. 
7. Let $R$ be a domain, $S$ a multiplicative set in $R$ not containing 0. Then we can view $S^{-1}R$ as a subring of the quotient field of $R$. Prove that

$$R = \bigcap_{\text{Max}(R)} m$$

where $\text{Max}(R)$ is the set of maximal ideals in $R$. 
CHAPTER VI

Domains

In this chapter, we study in greater depth special domains. The main goal is to study those domains that satisfy the analogue of the Fundamental Theorem of Arithmetic called Unique Factorization Domains or UFDs and explicit examples of these domains. The UFDs include PIDs and the historically interesting special class of PIDs that satisfy the division algorithm called euclidean domains. The explicit example of a euclidean domain, the Gaussian integers, is used to prove the classical two square theorems in Number Theory. We also prove Lagrange’s Four Square Theorem to indicate the method of proof developed by Fermat’s called Infinite Descent (a form of induction) to prove this theorem. We also introduce a more general class of commutative rings called Noetherian rings. These are the most important commutative rings in algebra and in algebraic geometry (as well as historically the most important because of the work of Hilbert). We shall study these rings in greater detail in a later chapter. For purposes here they represent a generalization of PIDs. From the viewpoint of induction, they represent rings where finite induction has a direct analogue. Theorems proven using Zorn’s Lemma over a commutative ring, usually do not need it if the ring is Noetherian.

29. Special Domains

We have seen that prime elements in \( \mathbb{Z} \) are those elements that satisfy one of two equivalent conditions: Let \( a, b \in \mathbb{Z} \). Then \( p \neq 0 \) is a prime if

1. whenever \( p = ab \), then either \( a = \pm 1 \) or \( b = \pm 1 \), i.e., either \( a \) or \( b \) is a unit or
2. whenever \( p \mid ab \) then \( p \mid a \) or \( p \mid b \).

The first was the original definition, the second Euclid’s Lemma. We have generalized the concept of a prime element in \( \mathbb{Z} \) to the concept of a prime ideal, but have not generalized the first equivalence in this way, nor will we, as its importance is related to elements. Elements satisfying the generalization of the first condition will be called irreducible elements. That elements in a domain factor into a product of
irreducible elements (up to units) holds for a large class of domains. What does not hold in general and what makes the two conditions not equivalent is an appropriate uniqueness statement. We shall carefully study domains in which we do have such a uniqueness statement. In general the second condition above is the more important as it gives more information about rings.

**Definition 29.1.** Let $R$ be a domain and $a$ a nonzero nonunit in $R$. We say $a$ is irreducible if whenever $a = bc$ with $b, c \in R$, then either $b$ or $c$ is a unit in $R$ (if $a$ is not irreducible, we say that it is reducible) and is a prime element if whenever $a \mid bc$ with $b, c \in R$, then $a \mid b$ or $a \mid c$.

Note that any associate of an irreducible element (respectively, prime element) is irreducible (respectively, prime).

**Remark 29.2.** Let $R$ be a domain and $p \in R$ a prime element. Then $p$ is irreducible i.e., being a prime element is stronger than being an irreducible element.

**Proof.** Suppose that $a = bc$ with $b, c \in R$. As $a$ is a prime element, either $a \mid b$ or $a \mid c$, say $a \mid b$. Then $b = ax$ some nonzero element $x \in R$. Therefore, we have $b = ax = bcx$ in the domain $R$, so cancellation implies that $1 = cx$, i.e., $c \in \mathbb{Z}^\times$ as needed.

We want to generalize the Fundamental Theorem of Arithmetic to a wider class of domains. Before, our approach relied on the concept of greatest common divisors. In general, our notion of a gcd of two elements below in a domain may not exist. Moreover, for integers, a gcd was always unique because we could make it positive as $\mathbb{Z}^\times = \{\pm1\}$. In general, we cannot do this, so must omit this condition that a gcd is unique if it exists.

**Definition 29.3.** Let $R$ be a domain, $a, b$ and $d$ nonzero elements in $R$. We call $d$ a greatest common divisor or gcd of $a$ and $b$ if:

(i) $d \mid a$ and $d \mid b$ in $R$.

(ii) If $c$ in $R$ satisfies $c \mid a$ and $c \mid b$ in $R$, then $c \mid d$ in $R$.

If a gcd of $a$ and $b$ exists and is a unit, then we say that $a$ and $b$ are relatively prime.

**Remarks 29.4.** Let $R$ be a domain and $a, b$ nonzero elements in $R$.

1. As mentioned above, a gcd of $a, b$ may not exist. For example, there exist elements in the subring $R := \{x + 2y\sqrt{-1} \mid x, y \in \mathbb{Z}\}$ of the Gaussian integers $\mathbb{Z}[\sqrt{-1}] := \{x + y\sqrt{-1} \mid x, y \in \mathbb{Z}\}$ not having a gcd.
2. If a gcd $d$ of $a, b$ exists it may not satisfy $d = ax + by$ for some $x, y \in R$, as it did for $\mathbb{Z}$.

3. If $d \in R$ satisfies $(d) = (a, b)$ then $d$ is a gcd of $a$ and $b$ (and, of course, there exist $x, y \in R$ satisfying $d = ax + by$).

4. If $d$ and $d'$ are two gcd’s of $a$ and $b$ in $R$, then $d \mid d'$ and $d' \mid d$, hence $d \approx d'$, i.e., $(d) = (d')$. Moreover, if this is the case, then any generator $c$ of $(d)$, i.e., an element $c$ in $R$ satisfying $(c) = (d)$, is a gcd of $a$ and $b$. In particular, a gcd of $a$ and $b$, if it exists, is unique up to units in $R$.

We codify some of the remarks above in the following:

Proposition 29.5. Let $R$ be a PID, $a$ and $b$ nonzero elements in $R$. Let $(d) = (a, b)$. Then $d$ is a gcd of $a$ and $b$ in $R$, unique up to units.

In particular, gcd’s of nonzero elements in a PID always exist and if $d'$ is a gcd of $a$ and $b$, then $d' = ax + by$ for some $x, y \in R$.

The proposition allows us to extend the analogue of Euclid’s Lemma to any PID.

Corollary 29.6. (Euclid’s Lemma) Let $R$ be a PID and $r \in R$. Then $r$ is irreducible if and only if $r$ is a prime element.

Proof. We have already shown that prime elements are irreducible. Conversely, assume that $r$ is an irreducible element in $R$ and satisfies $r \mid ab$ with $a, b \in R$. As $R$ is a PID, there exists a $c \in R$ satisfying $(c) = (r, a)$. Write $r = xc$ with $x \in R$. As $r$ is irreducible, either $c$ is a unit or $x$ is a unit in $R$. If $c$ is a unit, then $(c)$ is the unit ideal, hence $1 = ry + az$ in $R$ for some $y, z \in R$. It follows that $r \mid ryb + abz = b$ as needed. If $x$ is a unit then $(r) = (c) = (r, a)$, so $a \in (r)$, i.e., $r \mid a$. □

The Fundamental Theorem of Arithmetic relied on the equivalence of the notion of irreducible and prime elements. We wish to generalize this theorem to a wider class of domains, so we must say what we mean by uniqueness of factorization. For the integers this meant up to $\pm 1$, the units in $\mathbb{Z}$, so we will use the following generalization.

Definition 29.7. Let $R$ be a domain. We call $R$ a unique factorization domain or UFD if for any nonzero nonunit $r \in R$, we have

(i) $r = f_1 \cdots f_n$ for some irreducible elements $f_1, \ldots, f_n$ in $R$ for some $n \in \mathbb{Z}^+$,

(ii) If $r = f_1 \cdots f_n = g_1 \cdots g_m$ with $f_1, \ldots, f_n, g_1, \ldots, g_m$ irreducible elements in $R$ for some $n, m \in \mathbb{Z}^+$, then $n = m$ and there exists a permutation $\sigma \in S_n$ satisfying $f_i \approx g_{\sigma(i)}$ for all $i$. 

If \( r = uf_1 \cdots f_n \) with \( u \) a unit in \( R \) and \( f_i \in R \) for all \( i \), then we can ‘absorb’ the unit by replacing any one of the factors with an associate.

**Remark 29.8.** If \( R \) is a UFD, then, as desired, an element in \( R \) is prime if and only if it is irreducible.

**Proof.** We need only show if \( r \) is an irreducible element in \( R \) that it is a prime. So suppose that \( r \mid ab \) with \( a, b \in R \) nonunits. Write \( a = f_1 \cdots f_n \) and \( b = g_1 \cdots g_m \) with \( f_1, \ldots, f_n, g_1, \ldots, g_m \) irreducible elements in \( R \) for some \( n, m \in \mathbb{Z}^+ \). Since \( ab = rx \) for some \( x \in R \) by assumption, the uniqueness property in the definition of UFD’s implies \( r \approx f_i \) some \( i \) or \( r \approx g_j \) some \( j \), i.e., \( r \mid a \) or \( r \mid b \). \( \square \)

The uniqueness property for a UFD is really just that irreducible elements are prime:

**Proposition 29.9.** (Euclid’s Argument) Let \( R \) be a domain. Suppose that \( p_1 \cdots p_n = f_1 \cdots f_m \) with \( p_1, \ldots, p_n \) prime elements in \( R \) and \( f_1, \ldots, f_m \) irreducible elements in \( R \) for some \( n, m \in \mathbb{Z}^+ \). Then \( n = m \) and there exists a permutation \( \sigma \in S_n \) satisfying \( p_i \approx f_{\sigma(i)} \). In particular, each \( f_i \) is a prime element in \( R \).

**Proof.** As \( p_1 \mid f_1 \cdots f_m \) in \( R \), there exists an \( i \) such that \( f_i = up_1 \) for some element \( u \) in \( R \). Since \( f_i \) is irreducible and \( p_1 \) not a unit, \( u \) is a unit in \( R \). Now \( p_1 \cdots p_n = up_1 f_2 \cdots f_m \) in the domain \( R \). If \( n = 1 \) or \( m = 1 \) we must have \( m = n = 1 \) by cancellation, since the product of irreducibles is not a unit. So we may assume that \( m, n > 1 \). Cancellation shows \( p_2 \cdots p_n = uf_2 \cdots f_m \), and the result follows by induction. \( \square \)

An immediate consequence is the following:

**Corollary 29.10.** Let \( R \) be a domain. Suppose that every nonzero nonunit in \( R \) is a product of (finitely many) prime elements, i.e., \( a = p_1 \cdots p_n \), with \( p_1, \ldots, p_n \) prime elements in \( R \). Then \( R \) is a UFD.

The most important class of commutative rings is the following:

**Definition 29.11.** A commutative ring \( R \) is called a *Noetherian ring* if one of the following equivalent conditions hold:

(i) \( R \) satisfies the *ascending chain condition* or *ACC*, viz., given a (countable) chain of ideals

\[ \mathfrak{a}_1 \subset \mathfrak{a}_2 \subset \cdots \subset \mathfrak{a}_n \subset \cdots \]

in \( R \), there exists a positive integer \( N \) such that \( \mathfrak{a}_{N+i} = \mathfrak{a}_N \) for all \( i \geq 0 \).
We say that the chain stabilizes.
Equivalently, there exists no infinite chain of ideals in $R$, $\mathcal{B}_1 < \mathcal{B}_2 < \cdots < \mathcal{B}_n < \cdots$.

(ii) Every ideal $\mathfrak{A}$ in $R$ is finitely generated (fg), i.e., there exist $a_1, \ldots, a_n \in R$, some $n$, such that $\mathfrak{A} = (a_1, \ldots, a_n)$.

(iii) $R$ satisfies the Maximal Principle, i.e., any non-empty set of ideals $S$ in $R$ contains a maximal element relative to the relation $\subset$.

[Cf. this condition to the conclusion of Zorn’s Lemma.]

Remark 29.12. We leave it as an important exercise (cf. Exercise 29.22(6)) to show these three conditions are equivalent. However, the proof of this really needs the Axiom of Choice to prove ACC implies the Maximal Principle, and it is known that the Axiom of Choice is equivalent to Zorn’s Lemma. [Cf. Appendix A.] In fact, it is also known that the Axiom of Choice is also equivalent to the implication that ACC implies the Maximal Principle. For Noetherian rings one usually uses the Maximal Principle instead of directly using Zorn’s Lemma in proofs.

Using the second condition we have:

Examples 29.13. 1. Every PID is a Noetherian domain.

2. Let $R$ be a commutative ring. By Exercise 27.19(4), $R$ is Noetherian if and only if every prime ideal is finitely generated. Let $A = R[[t]]$. By Exercise 27.19(6), prime ideals in $A$ are finitely generated if prime ideals in $R$ are. [The converse is also true.] Hence if $R$ is Noetherian, so is $R[[t]]$.

Noetherian rings are important as they are the rings that arise in classical algebraic geometry. The collection of Noetherian rings also has the properties of being closed under basic ring constructions (and combinations of them): viz.,

(i) The homomorphic image of a Noetherian ring is Noetherian.

(ii) The localization of a Noetherian ring is Noetherian.

(iii) (Hilbert Basis Theorem) If $R$ is a Noetherian ring so is $R[t]$.

We shall leave the first two as exercises and prove the last in Section 40. (Cf. Theorem 40.1). It is not true, however, that a subring of a Noetherian ring is Noetherian. Indeed as any field is a Noetherian domain, one need only find a non-Noetherian domain to provide a counterexample. Can you do so?

We illustrate the importance of the Noetherian condition by the following interesting result.
Theorem 29.14. Let $R$ be a Noetherian domain and $r$ a nonzero non-unit in $R$. Then $r$ is a product of (finitely many) irreducible elements.

Proof. Let

$$S = \{ (a) \mid (0) < (a) < R \text{ with } a \text{ not a product of irreducible elements} \}.$$ 

Suppose that $S$ is non-empty, i.e., the result is false. Then there exists a principal ideal $(a) \in S$, a maximal element by the Maximal Principle. We call $(a)$ a maximal counterexample. Certainly, $a$ cannot be irreducible, so $a = bc$ for some non-units $b, c \in R$. But then $(a) < (b)$ and $(a) < (c)$, so by maximality, both $b$ and $c$ are products of irreducible elements, hence so is $a = bc$. □

The above proof is called proof by Noetherian induction. Note the similarity between this proof and that showing positive integers factor into primes. Further application of this type of decomposition into ‘smaller pieces’ for Noetherian rings is given in the exercises.

Theorem 29.15. Let $R$ be a PID. Then $R$ is a UFD.

Proof. As $R$ is a PID, it is a Noetherian domain. Hence every nonzero nonunit factors into a product of irreducible elements. As every irreducible element in a PID is a prime, factorization is (essentially) unique by Euclid’s Argument. □

Corollary 29.16. $\mathbb{Z}$ is a UFD.

Next we generalize the Division Algorithm, to study a special type of PID.

Definition 29.17. Let $R$ be a domain. Then $R$ is called a euclidean domain if there exists a function

$$\partial : R \setminus \{0\} \to \mathbb{Z}^+ \cup \{0\}$$

satisfying the Division Algorithm, i.e., for all $a, b \in R$ with $b$ nonzero, there exist $q, r \in R$ satisfying

(i) \quad $a = bq + r.$

(ii) \quad $r = 0$ or $\partial r < \partial b.$

We call $\partial a$ a euclidean function.

Remark 29.18. If we agree that $-\infty + n = -\infty$ for all integers $n$ and let $\partial 0 = -\infty$, we can write (ii) as

$$\partial r < \partial b.$$ 

[This is often done.]
Remark 29.19. Historically, a euclidian function was also made to satisfy

\[(iii) \quad \partial(ab) \geq \partial(a) \text{ for all nonzero } a, b \in R.\]

If \((iii)\) also holds, we shall call this a strong euclidean domain.

If \((iii)\) holds, then \(\partial(ab) > \partial(a)\) if \(b\) is not a unit and \(\partial(u) = \partial(1)\) for all units \(u \in R^\times\).

Theorem 29.20. If \(R\) is a euclidean domain, then \(R\) is a PID, hence a UFD.

**Proof.** Let \(0 < \mathfrak{A} \subset R\) be an ideal and

\[\emptyset \neq S = \{\partial a \mid 0 \neq a \in \mathfrak{A}\} \subset \mathbb{Z}^+ \cup \{0\}\]

(as \(\mathfrak{A} \neq 0\)). By the Well-ordering Principle, there exists an element \(a\) in \(\mathfrak{A}\) with \(\partial a \in S\) minimal.

**Claim.** \(\mathfrak{A} = (a)\) [Cf. the proof that \(\mathbb{Z}\) is a PID.]:

Let \(b \in \mathfrak{A}\), then \(b = aq + r\) in \(R\) with \(r = 0\) or \(\partial r < \partial a\). If \(r \neq 0\), then \(0 \neq r = b - aq\) lies in \(\mathfrak{A}\), contradicting the minimality of \(\partial a\) so \(r = 0\) and \(\mathfrak{A} = (a)\).

Remarks 29.21. (no justification)

1. \(\mathbb{Z}[t]\) is a UFD but not a PID.
2. \(\mathbb{Z}[\sqrt{-5}]\) is not a UFD.
3. \(\mathbb{Z}[\sqrt{-1}], \mathbb{Z}[\sqrt{-2}], \mathbb{Z}[\frac{-1-\sqrt{-3}}{2}], \mathbb{Z}[\frac{-1-\sqrt{-7}}{2}], \mathbb{Z}[\frac{-1-\sqrt{-11}}{2}]\) are (strong) euclidean domains.
4. \(\mathbb{Z}[\frac{-1-\sqrt{-19}}{2}]\) is a PID but not strong euclidean.
5. Any field is euclidean.
6. If \(F\) is a field, then the polynomial ring \(F[t]\) is euclidean with euclidean function \(\deg\).

Exercises 29.22.

1. Let \(R\) be a domain and \(a\) a nonzero nonunit in \(R\). Show that \(a\) is irreducible if and only if the principal ideal \((a)\) is maximal in the set \(\{(b) \mid b\) a nonzero nonunit in \(R\}\). In particular, if \(R\) is a PID, then every irreducible element in \(R\) is a prime element.
2. Produce elements \(a\) and \(b\) in the domain \(R := \{x+2y\sqrt{-1} \mid x, y \in \mathbb{Z}\}\) having no gcd. Prove your elements do not have a gcd.
3. Let \(\mathbb{Z} \subset R\) be a subring with \(R\) a UFD. Let \(d\) be the gcd of two nonzero integers \(a\) and \(b\) in \(\mathbb{Z}\). Show that \(d\) is still a gcd of \(a\) and \(b\) in \(R\).
4. Show 1 is a gcd for 2 and \( t \) in \( \mathbb{Z}[t] \), but there are no polynomials \( f, g \in \mathbb{Z}[t] \) satisfying \( 1 = 2f + tg \).

5. Let \( R \) be a UFD and \( a \) a nonzero element in \( R \). Show that the nilradical of \( R/(a) \) is the intersection of a finite number of prime ideals in \( R/(a) \), say \( \mathfrak{p}_1, \ldots, \mathfrak{p}_n \), and if \( \mathfrak{p} \) is any prime ideal in \( R/(a) \) then there exists an \( i \) such that \( \mathfrak{p}_i \subset \mathfrak{p} \). (Cf. Exercise 27.19 (9).)

6. Prove that the three conditions defining a Noetherian ring in (29.11) are indeed equivalent.

[Note: Technically to prove ACC implies the Maximal Principle, one needs the Axiom of Choice as mentioned in Remark 29.12. If you see where, you can invoke it.]

7. If in the previous exercise, you do not want to use the Axiom of Choice, but rather use Zorn’s Lemma, show that if all ideals of a commutative ring are finitely generated, then the Maximal Principle holds using Zorn’s Lemma.


9. A commutative ring \( R \) is called Artinian if it satisfies the following condition (called the descending chain condition or DCC): Any chain of ideals

\[
\mathfrak{A}_1 \supset \mathfrak{A}_2 \supset \cdots \supset \mathfrak{A}_n \supset \cdots
\]

(countable) in \( R \) stabilizes, i.e., there exists an integer \( N \) such that \( \mathfrak{A}_{N+i} = \mathfrak{A}_N \) for all \( i \geq 0 \). Equivalently, there exist no infinite chains

\[
\mathfrak{B}_1 > \mathfrak{B}_2 > \cdots > \mathfrak{B}_n > \cdots.
\]

Show that \( R \) is Artinian if and only if it satisfies the Minimal Principle which says that any non-empty collection of ideals in \( R \) has a minimal element (under set inclusion). (Cf. See the Note in Exercise (6) above. It applies here as well.)

10. If \( R \) is a domain, show that it is Artinian (cf. previous exercise) if and only if it is a field.

11. Let \( R \) be a Noetherian ring. Show if \( \varphi : R \to S \) is a ring epimorphism, then \( S \) is Noetherian.

12. If \( R \) is a Noetherian ring and \( S \) a multiplicative set in \( R \), then the localization \( S^{-1}R \) is a Noetherian ring. (Cf. Exercise 28.3(6).)

13. Let \( R \) be a Noetherian domain. Show that any nontrivial ideal of \( R \) contains a finite product of nonzero prime ideals, i.e., if \( 0 < \mathfrak{A} < R \) is an ideal, then there exist nonzero prime ideals \( \mathfrak{p}_1, \ldots, \mathfrak{p}_n \) in \( R \) such that \( \mathfrak{p}_1 \mathfrak{p}_2 \cdots \mathfrak{p}_n \subset \mathfrak{A} \).
14. Let $R$ be a Noetherian ring. Show that any ideal $\mathfrak{A} < R$ contains a finite product of prime ideals.

15. An ideal $\mathfrak{C}$ in a commutative ring $R$ is called irreducible if whenever $\mathfrak{C} = \mathfrak{A} \cap \mathfrak{B}$ for some ideals $\mathfrak{A}$ and $\mathfrak{B}$ in $R$, then either $\mathfrak{C} = \mathfrak{A}$ or $\mathfrak{C} = \mathfrak{B}$. Show if $R$ is Noetherian, then every ideal $\mathfrak{A} < R$ is a finite intersection of irreducible ideals of $R$, i.e., $\mathfrak{A} = \mathfrak{C}_1 \cap \cdots \cap \mathfrak{C}_n$, for some irreducible ideals $\mathfrak{C}_i$ in $R$.

16. Let $R$ be a Noetherian ring and $\mathfrak{A} < R$ an irreducible ideal (cf. Exercise 15 above). Show that $\mathfrak{A}$ is a primary ideal. (Cf. Exercise 27.19(10) for the definition of a primary ideal).

17. Let $f, g$ be polynomials with coefficients in a commutative ring $R$. Suppose the leading coefficient of $f$ is a unit (i.e., the coefficient of the highest degree term in $f$ is a unit). Show that there are polynomials $q$ and $r$ with coefficients in $R$ such that $g = fq + r$ with either $r = 0$ or the degree of $r$ is less than the degree of $f$. This says the Division Algorithm holds when dividing by a polynomial with unit leading term.

18. Let $F$ be a field. Show that $F[t]$, the ring of polynomials with coefficients in $F$, is a euclidean domain, hence a PID, hence a UFD.

19. Prove that a domain in which every prime ideal is principal is a PID.

20. Let $R$ be a strong euclidean domain, $a, b$ nonzero elements in $R$, Show that $b$ is a unit if and only if $\delta(b) = \delta(1)$ and if $b$ is a nonunit, then $\delta(ab) > \delta(a)$.

21. Let $R$ be a euclidean domain. Show directly, i.e., without using facts about PIDs or Noetherian domains, the following:

(a) Every irreducible element in $R$ is a prime element.

(b) If $R$ is a strong euclidean domain, then $R$ is a UFD.

22. Let $R$ be a UFD with quotient field $K$. Let $S$ be a multiplicative set in $R$ not containing zero. Show that the localization $S^{-1}R$ of $R$ at $S$ is a UFD. In particular, if $p$ is a prime ideal in $R$, then the localization $R_p := S^{-1}R$ of $R$ at $S = R \setminus p$,

$$R_p := \left\{ \frac{a}{b} \mid a, b \in R, \ b \notin p \right\} \subset K$$

is also a UFD.

30. Characterization of UFDs

We prove a generalization of our result that a PID is a UFD by giving Kaplansky’s characterization of UFD’s.
Theorem 30.1. (Kaplansky) Let $R$ be a domain. Then $R$ is a UFD if and only if every nonzero prime ideal contains a prime element.

Proof. ($\Rightarrow$): Let $0 < \mathfrak{p} < R$ be a prime ideal and $a$ a nonzero element in $\mathfrak{p}$. Then $a = f_1 \cdots f_r$, with $f_1, \ldots, f_r$ irreducible in $R$. Since $\mathfrak{p}$ is a prime ideal, there exist an $i$ such that $f_i \in \mathfrak{p}$. As $R$ is a UFD, $f_i$ is a prime element.

($\Leftarrow$): Let $S = \{ r \in R \mid r \neq 0 \text{ and } r \text{ is a unit or a product of prime elements} \}$. Clearly, $S \neq \emptyset$ is a multiplicative set and $(0)$ excludes $S$. If we show $S = R \setminus \{0\}$, then $R$ is a UFD by Euclid’s Argument, i.e., we shall be done.

Claim. $S$ is a saturated multiplicative set, i.e., if $a, b \in R$, then

$$a, b \in S \text{ if and only if } ab \in S,$$

As $S$ is a multiplicative set, we need only show if $ab \in S$, then $a$ and $b$ lie in $S$. Note that $R^\times$ is a saturated multiplicative set, and by definition a subset of $S$, so we may assume that neither $a$ nor $b$ is a unit. Write $ab = p_1 \cdots p_r$ with $p_1, \ldots, p_r$ primes hence each $p_i$ lies in $S$. As $p_1 | ab$ and $p_1$ is a prime element, $p_1 | a$ or $p_1 | b$, say $p_1 | a$. Write $a = p_1a_1$ with $a_1 \in R$, hence $p_1a_1b = p_1p_2 \cdots p_r$ in the domain $R$. By cancellation, we conclude that $a_1b = p_2 \cdots p_r$ in $S$. By induction on $r$, we have $a_1, b \in S$. As $S$ is a multiplicative set, $a \in S$. This proves the claim.

As $(0)$ excludes $S$, by Krull’s Theorem, there exists a prime ideal $\mathfrak{p} \in R$ excluding $S$ and maximal with respect to excluding $S$. If $0 < \mathfrak{p}$, then there exists a prime element $p \in \mathfrak{p}$, hence $p \in S \cap \mathfrak{p}$, a contradiction. Consequently, $(0) = \mathfrak{p}$ is maximal with respect to excluding $S$. In particular, if $a \in R$ is nonzero and satisfies $(0) < (a)$, we have $(a) \cap S \neq \emptyset$. We conclude that there exists an $r \in R$ such that $ar \in S$. As $S$ is saturated by the claim, $a \in S$ as needed.

Remarks 30.2. 1. If $R$ is a Noetherian ring, then, a priori, we do not need Zorn’s Lemma to prove Krull’s Theorem, hence in the proof above, if $R$ is a Noetherian domain. However, as mentioned in Remark 29.12 this is really illusory.

2. We shall see that $\mathbb{C}[t_1, \ldots, t_n]$ is a UFD, Using this one can show that the geometric meaning of Kaplansky’s Theorem in the case of $\mathbb{C}^n$ is: Any hypersurface in $\mathbb{C}^n$ given by polynomial equations (in several variables) satisfying an irreducible condition can be defined by a single irreducible polynomial.
31. Gaussian Integers

In this section, we study the Gaussian integers \( \mathbb{Z}[\sqrt{-1}] \), a domain as it is a subring of \( \mathbb{C} \). The map \( \overline{\cdot} : \mathbb{C} \to \mathbb{C} \) given by \( x + y\sqrt{-1} \mapsto x - y\sqrt{-1} \), \( x,y \in \mathbb{R} \), will denote complex conjugation. (As \( \mathbb{C} \) is a real vector space on \( \{1, \sqrt{-1}\} \), any element \( z \) in \( \mathbb{C} \) can be written \( z = x + y\sqrt{-1} \) for some unique real numbers \( x \) and \( y \). This is a field automorphism that fixes the reals, i.e., is the identity on the reals. In fact, \( \overline{z} = z \) if and only if \( z \in \mathbb{R} \). Clearly, it restricts to a ring automorphism of the Gaussian integers. Let \( N : \mathbb{C} \to \mathbb{R}^+ \cup \{0\} \) be the map defined by \( z \mapsto z\overline{z} \). It is called the norm map.

Remarks 31.1. 1. If \( R \) is a ring, we let \( \text{Aut}_{\text{ring}}(R) \) denote the set of ring automorphisms of \( R \). It is a group under composition. As \( \overline{\overline{z}} = z \) for all \( z \in \mathbb{C} \), complex conjugation has order two in \( \text{Aut}_{\text{ring}}(\mathbb{C}) \). A ring automorphism of commutative rings of order one or two is called an involution. So complex conjugation is an involution on \( \mathbb{C} \) and on \( \mathbb{Z}[\sqrt{-1}] \).

2. Let \( x, y \in \mathbb{R} \), then \( N(x + y\sqrt{-1}) = x^2 + y^2 \), which is positive unless \( x = 0 = y \). This indicates that the norm map should be of interest when studying sums of two squares.

3. \( N : (\mathbb{C}^\times, \cdot) \to (\mathbb{R}^\times, \cdot) \) is a monoid homomorphism, i.e., \( N(1) = 1 \) and it is multiplicative:

\[
N(z_1z_2) = N(z_1)N(z_2) \quad \text{for all} \quad z_1, z_2 \in \mathbb{C}^\times.
\]

In particular, if \( x_1, x_2, y_1, y_2 \in \mathbb{C} \), we have the two square identity:

\[
(x_1^2 + y_1^2)(x_2^2 + y_2^2) = (x_1x_2 - y_1y_2)^2 + (x_1y_2 + x_2y_1)^2,
\]
the product of sums of two squares is a sum of two squares. Of course, once we have this identity, we see that this holds in any commutative ring by multiplying out.

4. \( N(z) = N(\overline{z}) \) for all \( z \in \mathbb{C} \).

**Lemma 31.2.** \( \mathbb{Z}[\sqrt{-1}] \) is a domain and satisfies the following:

1. \( N|_{\mathbb{Z}[\sqrt{-1}]} : \mathbb{Z}[\sqrt{-1}] \to \mathbb{Z}^+ \cup \{0\} \) and satisfies \( N|_{\mathbb{Z}[\sqrt{-1}]}(\alpha) = 0 \) if and only if \( \alpha = 0 \). We shall write \( N \) for \( N|_{\mathbb{Z}[\sqrt{-1}]} \).

2. \( \mathbb{Z}[\sqrt{-1}]^\times = \{ \pm 1, \pm \sqrt{-1} \} \).

3. Let \( a, b \in \mathbb{Z} \). Then \( a + b\sqrt{-1} \), is a root of the monic polynomial (i.e., leading coefficient is one), \( t^2 - 2at + (a^2 + b^2) \in \mathbb{Z}[t] \).

**Proof.** We only prove (2) leaving the rest as an easy check. If \( uv = 1 \) with \( u,v \in \mathbb{Z}[t] \), then \( 1 = N(1) = N(u)N(v) \), so we have \( N(u) \in \mathbb{Z}^\times \cap \mathbb{Z}^+ = \{ 1 \} \). As \( x^2 + y^2 = 1 \) with \( x,y \in \mathbb{Z} \) if and only if \( x = \pm 1, y = 0 \) or \( x = 0, y = \pm 1 \), we conclude that

\[
\mathbb{Z}[\sqrt{-1}]^\times \subset \{ \beta \mid N(\beta) = 1 \} \subset \{ \pm 1, \pm \sqrt{-1} \}.
\]

Clearly, \( \{ \pm 1, \pm \sqrt{-1} \} \subset \mathbb{Z}[\sqrt{-1}]^\times \), establishing the lemma. \( \square \)

**Theorem 31.3.** The domain \( \mathbb{Z}[\sqrt{-1}] \) is a (strong) euclidean domain, hence a PID and so a UFD.

**Proof. Claim.** \( N \) is a euclidean function on \( \mathbb{Z}[\sqrt{-1}] \). (Consequently, \( \mathbb{Z}[\sqrt{-1}] \) is a strong euclidean domain.):

Let \( \alpha \) and \( \beta \) be elements in \( \mathbb{Z}[\sqrt{-1}] \) with \( \beta \) nonzero. We must show that there exist a \( \gamma \) and \( \rho \) in \( \mathbb{Z}[\sqrt{-1}] \) satisfying \( \alpha = \beta \gamma + \rho \) with \( \rho = 0 \) or \( N(\rho) < N(\beta) \).

In \( \mathbb{C} \), we can write

\[
\frac{\alpha}{\beta} = \frac{\alpha \overline{\beta}}{\beta \overline{\beta}} = \frac{\alpha \overline{\beta}}{N(\beta)} = r + s \sqrt{-1}
\]

for some \( r, s \in \mathbb{Q} \) (as \( N(\beta) \in \mathbb{Z} \)). Choose \( m, n \in \mathbb{Z} \) satisfying

\[
|r - m| \leq 1/2 \quad \text{and} \quad |s - n| \leq 1/2.
\]

Set \( \gamma = m + n\sqrt{-1} \) and \( \rho = \alpha - \beta \gamma \) in \( \mathbb{Z}[\sqrt{-1}] \). If \( \rho = 0 \), we are done, so suppose not. As \( N : \mathbb{C}^\times \to \mathbb{R}^+ \) and is multiplicative, we have

\[
N(\rho) = N(\alpha - \beta \gamma) = N(\beta(\frac{\alpha}{\beta} - \gamma)) = N(\beta)N(\frac{\alpha}{\beta} - \gamma)
= N(\beta)N((r - m) + (s - n)\sqrt{-1}) = N(\beta)[(r - m)^2 + (s - n)^2]
\leq N(\beta)(\frac{1}{4} + \frac{1}{4}) < N(\beta).
\]
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Note in the last line of the proof, our estimate on \( N(\rho) \) was \( \leq \frac{1}{2} N(\beta) \), so we had some room to work. This explains why the domains in Remark 29.21(3) work — compute the restriction of \( N \) to each of these domains.

As mentioned before, we can use the Gaussian integers to study sums of squares of integers. In particular, \( m = x^2 + y^2 \) with \( x, y \) integers if and only if \( m = N(x + y\sqrt{-1}) \). We wish to determine all the non-negative integers that are sums of two squares. We have seen that a product of two integers each a sum of two squares of integers is itself a sum of two squares. By the Fundamental Theorem of Arithmetic, each non-negative integer factors as \( p_1^{e_1} \cdots p_r^{e_r} \) with the \( p_i \) distinct positive primes and the \( e_i \geq 1 \). If \( e_i \) is even then \( p_i^{e_i} \) is a square (so a sum of two squares). Hence we are reduced to determining which products of distinct primes are sums of two squares of integers. In particular, we must determine which primes are sums of two squares of integers. As 2 is a sum of two squares, we are reduced to odd primes. If \( x \in \mathbb{Z} \), then \( x^2 \equiv 0, 1 \) mod 4. In particular, a sum of two integer squares can only be congruent to 0, 1 or 2 modulo four. In particular, no integer \( n \) satisfying \( n \equiv 3 \) mod 4 can be a sum of two squares of integers, e.g., 3 and 7 are not sums of two squares.

The key to the solution is the following:

**Lemma 31.4.** Let \( p \) be a positive prime in \( \mathbb{Z} \) satisfying \( p \equiv 1 \) mod 4. Then there exists an integer \( x \) satisfying \( p \mid x^2 + 1 \) in \( \mathbb{Z} \). In particular, \(-1\) is a square modulo \( p \) if \( p \equiv 1 \) mod 4.

**Proof.** Let \( x = 1 \cdot 2 \cdot 3 \cdots \frac{p-1}{2} \) in \( \mathbb{Z} \). By Wilson’s Theorem (Exercise 10.16(13)), we have

\[
-1 \equiv (p - 1)! = (1 \cdot 2 \cdot 3 \cdots \frac{p-1}{2})(p - \frac{p-1}{2} \cdots p - 3 \cdot p - 2 \cdot p - 1)
\equiv (1 \cdot 2 \cdot 3 \cdots \frac{p-1}{2})(- \frac{p-1}{2} \cdots - 3 \cdot -2 \cdot -1) = (-1)^{\frac{p-1}{2}} x^2
\equiv x^2 \mod p,
\]

so \( p \mid x^2 + 1 \). \( \square \)

**Theorem 31.5.** (Fermat) Let \( p \) be a positive prime congruent to 1 modulo 4. Then there exist integers \( x \) and \( y \) such that \( p = x^2 + y^2 \).

**Proof.** By the lemma there exists an integer \( x \) such that \( p \mid x^2 + 1 \). Hence \( p \mid (x + \sqrt{-1})(x - \sqrt{-1}) \) in the UFD \( \mathbb{Z}[\sqrt{-1}] \).

**Claim.** \( p \) is not irreducible (i.e., it is reducible) in \( \mathbb{Z}[\sqrt{-1}] \):
Suppose that \( p \) is irreducible in \( \mathbb{Z}[\sqrt{-1}] \). Then as \( \mathbb{Z}[\sqrt{-1}] \) is a UFD, it is a prime. Hence \( p \mid x + \sqrt{-1} \) or \( p \mid x - \sqrt{-1} \). If \( p \mid x + \sqrt{-1} \),
then \( p(a + b\sqrt{-1}) = x + \sqrt{-1} \) for some integers \( a \) and \( b \). It follows that \( pb = 1 \) in \( \mathbb{Z} \) (why?), which is impossible. [Or if \( p \mid x + \sqrt{-1} \), then \( p = p \mid x - \sqrt{-1} \) also. It follows that \( p \mid 2x \) hence \( p \mid x \) as \( p \) is odd.]

Similarly, \( p \nmid x - \sqrt{-1} \), so \( p \) is not irreducible in \( \mathbb{Z}[\sqrt{-1}] \), proving the claim.

Using the claim, we can write \( p = \pi \alpha \) for some \( \pi, \alpha \in \mathbb{Z}[\sqrt{-1}] \), with \( \pi \) irreducible in \( \mathbb{Z}[\sqrt{-1}] \) and \( \alpha \) not a unit in \( \mathbb{Z}[\sqrt{-1}] \). Taking the norm of this equation yields

\[
p^2 = N(p) = N(\pi)N(\alpha)
\]

in \( \mathbb{Z} \). As \( N(\pi) \) and \( N(\alpha) \) lie in \( \mathbb{Z}^+ \) with \( \mathbb{Z} \) a UFD, we must have

\[
\pi \bar{\pi} = N(\pi) = p = N(\alpha)
\]

is a sum of two squares, proving the theorem. In addition, we have also shown, in view of this last equation, that \( \bar{\pi} = \alpha \) is also irreducible and \( p = \pi \bar{\pi} \) is a factorization of \( p \) in \( \mathbb{Z}[\sqrt{-1}] \). \( \square \)

To categorize those integers that are sums of two squares, we need the following lemma:

**Lemma 31.6.** Let \( a \) and \( b \) be two nonzero relatively prime integers and \( p \) a (positive) odd prime such that \( p \mid a^2 + b^2 \). Then \( p \equiv 1 \mod 4 \).

**Proof.** As \( p \mid a^2 + b^2 \), we have \( p \mid a \) if and only if \( p \mid b \). Since \( a \) and \( b \) are relatively prime, \( p \) cannot divide \( a \) or \( b \). By Fermat’s Little Theorem 10.14, we know that \( b^{p-3}b^2 = b^{p-1} \equiv 1 \mod p \). Consequently, \( a^2 + b^2 \equiv 0 \mod p \), i.e., \( a^2 \equiv -b^2 \mod p \), and \( p - 3 \) even imply that

\[
(b^{p-3}a)^2 = b^{p-3}a^2 \equiv -b^{p-3}b^2 \equiv -1 \mod p.
\]

Let \( x = b^{p-3}a \). Then \( x^2 \equiv -1 \mod p \) which means that

\[
1 \equiv x^{p-1} = (x^2)^{\frac{p-1}{2}} \equiv (-1)^{\frac{p-1}{2}} \mod p.
\]

If \( p \equiv 3 \mod 4 \), then we conclude that \( 1 \equiv -1 \mod p \), which is impossible. \( \square \)

The next result solves our stated problem. We leave its proof as an exercise (that you should do).

**Theorem 31.7.** Let \( n = p_1^{e_1} \cdots p_r^{e_r} \) in \( \mathbb{Z} \) with \( 0 < p_1 < \cdots < p_r \) primes, and \( e_i > 0 \) for \( i = 1, \ldots, r \). Then \( n \) is a sum of two squares in \( \mathbb{Z} \) if and only if \( e_i \) are even integers whenever \( p_i \equiv 3 \mod 4 \) for \( i = 1, \ldots, r \).

Let \( p \) be a (positive) odd prime. By Lemma 31.6, we know that

\[
-1 \text{ is } \begin{cases} 
\text{a square mod } p & \text{if } p \equiv 1 \mod 4, \\
\text{not a square mod } p & \text{if } p \equiv 3 \mod 4.
\end{cases}
\]
Definition 31.8. If $a$ is an integer not divisible by $p$, define the Legendre symbol
\[
\left( \frac{a}{p} \right) := \begin{cases} 
+1 & \text{if } a \mod p \text{ is a square.} \\
-1 & \text{if } a \mod p \text{ is not a square.}
\end{cases}
\]
It is also convenient to define $(\frac{b}{p}) = 0$ if $p | b$.

Example 31.9. $(\frac{0}{7}) = 0$, $(\frac{1}{7}) = (\frac{2}{7}) = (\frac{4}{7}) = 1$, and $(\frac{3}{7}) = (\frac{5}{7}) = (\frac{6}{7}) = -1$

We have proven:

Proposition 31.10. (Euler’s Formula) Let $p$ be a (positive) odd prime. Then
\[
\left( \frac{-1}{p} \right) = (-1)^{\frac{p-1}{2}} = \begin{cases} 
+1 & \text{if } p \equiv 1 \mod 4. \\
-1 & \text{if } p \equiv 3 \mod 4.
\end{cases}
\]

Corollary 31.11. The set $\{4n + 1 \mid n \in \mathbb{Z}\}$ contains infinitely many primes.

Proof. We prove the stronger result that the set $\{8n + 5 \mid n \in \mathbb{Z}\}$ contains infinitely many primes.

Check. If $n \in \mathbb{Z}$, then $n^2 \equiv 0, 1, 4 \mod 8$. In particular, if $n$ is odd, then $n^2 \equiv 1 \mod 8$.

Let $p$ be an odd prime and set $N = (3 \cdot 5 \cdot 7 \cdots p)^2 + 2^2 \equiv 5 \mod 8$. If $q$ is an odd prime such that $q | N$, then by Lemma 31.6, we have $q \equiv 1 \mod 4$, hence $q \equiv 1 \mod 8$ or $q \equiv 5 \mod 8$. If $N$ is a product of primes all congruent to 1 modulo 8, then $N \equiv 1 \mod 8$, a contradiction. Therefore, there exists a prime $q$ satisfying $q \equiv 5 \mod 8$ and $q | N$. As none of the primes $3, 5, 7, \ldots, p$ divides $N$, we must have $q > p$, so there exists a prime larger that $p$ congruent to 5 modulo 8.

Remark 31.12. It is still an open question whether the set of integers $\{x^2 + 1 \mid x \in \mathbb{Z}\}$ contains infinitely many primes.

The following is a well-known theorem, although we shall not prove it. The usual proof uses complex analysis, although there is now an elementary (but tricky) proof that does not.

Theorem 31.13. (Dirichlet’s Theorem on Primes in an Arithmetic Progression) Let $a$ and $b$ be two nonzero relatively prime integers. Then there exist infinitely many primes $p$ satisfying $p \equiv a \mod b$. 

We shall prove that this is true for $a = 1$ in Proposition 58.10 below by elementary means. Other questions arise about arithmetic progressions, i.e., equally spaced integers, e.g., $a, a+b, a+2b, a+3b, \ldots$ with $a$ and $b$ integers. The distance between two consecutive integers in an arithmetic progression is called the spacing of the progression, e.g., $b$ is the spacing in the example. Such a progression can be finite or infinite. The most spectacular theorem is the following proven in 2004:

**Theorem 31.14.** (Green-Tao) Let $N$ be a positive integer. Then there exists an arithmetic progression consisting of $N$ primes.

The theorem shows that for any $N$ there exists infinitely many arithmetic progressions consisting of $N$ primes, but it does not indicate anything about the spacing. Turning this theorem around, one can ask does there exist infinitely many primes in an arithmetic progression with a given spacing. The most famous example of this is

**Question 31.15.** (Twin Prime Conjecture) Do there exist infinitely many twin prime, i.e., primes pairs of the form $p, p+2$.

Euler had shown that $\sum_{p \text{ prime}} 1/p$ was infinite, so there was some hope that $\sum_{p \text{ twin prime}} 1/p$ was also. Brun showed this to be false in 1915.

We shall also study the Legendre symbol further proving the following wonderful (albeit seemingly strange) theorem of Gauss upon which you should meditate (and we shall prove it in Theorem 58.18 below):

**Theorem 31.16.** (Quadratic Reciprocity) Let $p$ and $q$ be distinct odd primes. Then

$$\left(\frac{p}{q}\right)\left(\frac{q}{p}\right) = (-1)^{\frac{p-1}{2}\frac{q-1}{2}}.$$  

**Exercises 31.17.**

1. Show that the evaluation map $e_{\sqrt{-1}} : \mathbb{Z}[t] \rightarrow \mathbb{C}$ defined by $f \mapsto f(\sqrt{-1})$ is a ring homomorphism with kernel $(t^2 + 1)$ and image the Gaussian integers. [Hint: Use Exercise 29.22(17).]

2. Let $R = \mathbb{Z}[\sqrt{-1}]$ and $n = p_1^{e_1} \cdots p_r^{e_r}$ be the standard factorization of the integer $n > 1$. Show that the following are equivalent:

   (i) $n$ is a sum of two squares.
   (ii) $n = N(\alpha)$ for some $\alpha \in R$.
   (iii) If $p_i \equiv 3 \pmod{4}$, then $e_i$ is even. 

   [Hint: Use Lemma 31.6.]
3. Show the following
   (i) Let $\alpha$ be an element in $\mathbb{Z}[\sqrt{-1}]$ such that $N(\alpha)$ is a prime or the square of a prime in $\mathbb{Z}$. Then $\alpha$ is a prime element or a product of two prime elements in $\mathbb{Z}[\sqrt{-1}]$.
   (ii) Let $\pi$ be an prime element in $\mathbb{Z}[\sqrt{-1}]$. Show $N(\pi)$ is a prime or the square of a prime in $\mathbb{Z}$.

4. Determine all prime elements, up to units, in $\mathbb{Z}[\sqrt{-1}]$.

5. Show that $\mathbb{Z}[\sqrt{-2}]$ is a (strong) euclidean domain.

6. Prove that $\mathbb{Z}\left[-\frac{1 - \sqrt{3}}{2}\right]$ and $\mathbb{Z}[\sqrt{3}]$ are (strong) euclidean domains.

7. Prove that $\mathbb{Z}\left[-\frac{1 - \sqrt{-d}}{2}\right]$ is a (strong) euclidean domain for $d = 7$ and $d = 11$.

8. Let $R = \mathbb{Z}[\sqrt{-d}] = \{a + b\sqrt{-d} \mid a, b \in \mathbb{Z}\}$, a subring of $\mathbb{C}$ with $d$ a positive square-free integer. Let $N : R \to \mathbb{Z}$ be the norm map, so $\alpha = a + b\sqrt{-d} \mapsto \alpha \bar{\alpha} = a^2 + db^2$. Show all of the following:
   (i) The field of quotients of $R$ is $\mathbb{Q}[\sqrt{-d}] = \{a + b\sqrt{-d} \mid a, b \in \mathbb{Q}\}$.
   (ii) $N : R \setminus \{0\} \to \mathbb{Z}$ is a monoid homomorphism.
   (iii) $R^\times = \{\alpha \in R \mid N(\alpha) = 1\}$ and compute this group for all $d$.
   (iv) The element $\alpha$ is irreducible in $R$ if $N(\alpha)$ is a prime. Is the converse true?
   (v) Suppose $d \geq 3$, then $2$ is irreducible but not prime in $R$.

9. Let $R = \mathbb{Z}[\sqrt{-5}]$. Show the following:
   (i) The elements $2, 3, 1 + \sqrt{-5}$, and $1 - \sqrt{-5}$ are all irreducible but no two are associates.
   (ii) None of the elements $2, 3, 1 + \sqrt{-5}$, and $1 - \sqrt{-5}$ are prime.

10. Let $R = \mathbb{Z}[\sqrt{-5}]$. Let $\mathfrak{P} = (2, 1 + \sqrt{-5})$. Show
    (i) $\mathfrak{P}^2 = (2)$ in $R$.
    (ii) $\mathfrak{P}$ is a maximal ideal.
    (iii) $\mathfrak{P}$ is not a principal ideal.

32. Addendum: The Four Square Theorem

In the previous section, we characterized those positive elements that are sums of two squares. In this section, we prove Lagrange’s theorem that every positive integer is a sum of four squares. The proof is based on a computation of Euler that shows a product of two integers each a sum four integer squares is a sum of four integer squares.
Although he did not know it at the time, this formula comes from the Hamiltonian quaternions, the first known division ring that was not commutative. We begin by constructing it.

**Construction 32.1.** Let $H$ be a four dimensional real vector space with basis $\{1, i, j, k\}$. We view 1 as $R$, so $R \subset H$. We make $H$ into a ring by defining a multiplication on this basis and extend linearly to the whole space. The multiplication of the basis elements is given by:

$$i^2 = -1, \quad j^2 = -1, \quad ij = -ji = k.$$ 

Extending this linearly give a multiplication on $H$ as follows:

$$(x_0 1 + x_1 i + x_2 j + x_3 k) \cdot (y_0 1 + y_1 i + y_2 j + y_3 k)$$

$$= (x_0 y_0 - x_1 y_1 - x_2 y_2 - x_3 y_3)1 + (x_0 y_1 + x_1 y_0 + x_2 y_3 - x_3 y_2)i$$

$$+ (x_0 y_2 + x_2 y_0 + x_3 y_1 - x_1 y_3)j + (x_0 y_3 + x_3 y_0 + x_1 y_2 - x_2 y_1)k.$$ 

By a straight-forward (but arduous) computation, one now checks that this makes $H$ into a ring with $1_H = 1_R$. Note that the subring generated by $\{1, i\}$ is isomorphic to $C$ (as are the subrings generated by $\{1, j\}$ and $\{1, k\}$, respectively). The ring $H$ is not commutative as its center, i.e., the subring $\{x \in H \mid xy = yx \text{ for all } y \in H\}$ of $H$, is precisely $R$. [Why is the center a subring?] Note also that $H^\times$ contains the *quaternion group* $\{1, i, j, ij, -1, -i, -j, -ji\}$.

Analogous to $C$, there is a *quaternion conjugation* map:

$$- : H \to H$$

given by $x_0 1 + x_1 i + x_2 j + x_3 k := x_0 1 - x_1 i - x_2 j - x_3 k$.

It is easy to check that for all $x, y$ in $H$, we have

$$\overline{1} = 1$$

$$\overline{x + y} = \overline{x} + \overline{y}$$

$$\overline{xy} = \overline{y} \overline{x}.$$ 

Therefore, the quaternion conjugation map satisfies the properties of a ring homomorphism except that it reverses multiplication. Such a map of rings is called a *ring antihomomorphism*. Clearly, this map is also a surjective linear transformation of a real vector space. Hence it is bijective, since $H$ is a finite dimensional vector space over $R$. Consequently, the map is, in fact, a *ring antiautomorphism*. (Note the composition of two ring antihomomorphisms is a ring homomorphism.) In addition, for all $x$ in $H$, we have

$$\overline{\overline{x}} = x,$$

so $-$ is an *involution*, i.e., an antiautomorphism satisfying the composition $- \circ -$ is the identity. (Cf. with complex conjugation.) We also
have a norm map:

\[ N : \mathcal{H} \to \mathbb{R}^+ \cup \{0\} \]

that is checked to be multiplicative, i.e., \( N(xy) = N(x)N(y) \) for all \( x, y \) in \( \mathcal{H} \). [Note that if we identify \( \mathbb{C} \) as the subring of \( \mathcal{H} \) generated by \( \{1, i\} \), then \( N|_{\mathbb{C}} \) is the usual norm map on \( \mathbb{C} \).] This is useful as the norm of \( x_01 + x_1i + x_2j + x_3k \) is

\[ N(x_01 + x_1i + x_2j + x_3k) = x_0^2 + x_1^2 + x_2^2 + x_3^2, \]

a sum of four squares. The multiplicativity of this norm map yields Euler’s Equation: the four square identity

\[
(x_0^2 + x_1^2 + x_2^2 + x_3^2) \cdot (y_0^2 + y_1^2 + y_2^2 + y_3^2) \\
= (x_0y_0 - x_1y_1 - x_2y_2 - x_3y_3)^2 + (x_0y_1 + x_1y_0 + x_2y_3 - x_3y_2)^2 \\
+ (x_0y_2 + x_2y_0 + x_3y_1 - x_1y_3)^2 + (x_0y_3 + x_3y_0 + x_1y_2 - x_2y_1)^2,
\]

the formula that a product of two sums of fours squares is a sum of four squares, generalizing the norm on the complex numbers giving a formula for the product of sums of two squares. This is the equation that we shall need to prove Lagrange’s Theorem that any positive integer is a sum of four squares. Of course, we do not need the quaternions to establish Euler’s Equation — we need only multiply the right hand side out to see that it holds in any commutative ring. Indeed Euler found this equation without knowing about quaternions. It does, however, explain why such an equation exists. It can also be shown that there is a formula that a product of two sums of eight squares is a sum of eight squares. It arises from a generalization of the quaternions called the octonians, a algebraic structure that looks a division ring but does not satisfy the associative law for multiplication. So the two square identity arises from fields, the four square identity from the quaternions (tossed out commutativity), the eight square identity from the octonians (toss out commutativity and associativity). Is there a sixteen square identity arising from a generalization of these algebraic structures by throwing out some other property of rings? Hurwitz showed not.

As with the norm map on \( \mathbb{C} \), it follows that \( N(z) = 0 \) if and only if \( z = 0 \). In particular, if \( z \) is nonzero, then \( z\overline{z}/N(z) = 1 = \overline{z}z/N(z) \), so \( z \) has a multiplicative inverse. Therefore, \( \mathcal{H} \) is a division ring called the Hamiltonian quaternions. By definition, \( ij \neq ji \), so \( \mathcal{H} \) is not a field.

**Remark 32.2.** Let \( V \) be a finite dimensional real vector space that is also a division ring and contains \( \mathbb{R} \) in its center. Then it can be shown that \( V \) is ring isomorphic to \( \mathbb{R}, \mathbb{C}, \) or \( \mathcal{H} \). In particular, \( \dim_{\mathbb{R}} V = 1, 2, \) or 4. We shall show this in Section 101.
To show that every positive integer is a sum of four squares, we shall use the following:

**Lemma 32.3.** \(1 \leq m < p\) and \(mp\) is a sum of four integer squares.

**Proof.** Let

\[ S_1 := \{0^2, 1^2, \ldots, \frac{(p-1)^2}{2}\} \]

\[ S_2 := \{-0^2 - 1, -1^2 - 1, \ldots, -\frac{(p-1)^2}{2} - 1\}. \]

If \(x^2 = y^2 \mod p\) with \(x, y\) integers, then \(p \mid x^2 - y^2 = (x + y)(x - y)\).
Consequently, \(p \mid x + y\) or \(p \mid x - y\). It follows (check) that no two elements in \(S_1\) are congruent modulo \(p\) and no two elements in \(S_2\) are congruent modulo \(p\). Since

\[ |S_1| + |S_2| = \left(\frac{p-1}{2} + 1\right) + \left(\frac{p-1}{2} + 1\right) = p + 1, \]

there exist \(x^2 \in S_1\) and \(-y^2 - 1 \in S_2\) satisfying \(0 \neq x^2 + y^2 + 1\) and

\[ x^2 + y^2 + 1 \equiv 0 \mod p. \]

Hence there exists a positive integer \(m\) such that

\[ mp = x^2 + y^2 + 1 \leq \left(\frac{p-1}{2}\right)^2 + \left(\frac{p-1}{2}\right)^2 + 1 < \frac{p^2}{2} + 1 < p^2. \]

So \(1 \leq m < p\). \(\square\)

**Remarks 32.4.**

1. Of course, the proof shows that the element \(mp\) is a sum of three squares.

2. The same proof shows that \(-1\) is a sum of two squares in \(\mathbb{Z}/p\mathbb{Z}\) for any prime \(p\). In fact, any element \(a\) in \(\mathbb{Z}/p\mathbb{Z}\) is a sum of two squares — replace the \(-1\) in the proof by \(a\) in \(S_1\).

3. A similar (but not obvious) proof works to show that any element in a finite field is a sum of two squares.
   [Note that if char\((F) = 2\), then the map \(F \to F\) given by \(x \mapsto x^2\) is a (field) monomorphism, so an automorphism if \(F\) is a finite field.]

Historically, one of the main forms of induction used to prove rather deep theorems was the idea of Fermat, called *Fermat descent*. We prove our theorem about four squares using this method.

**Theorem 32.5.** (Lagrange). Every positive integer is a sum of four squares of integers.
Proof. We know that $2 = 1 + 1 + 0 + 0$ is a sum of four squares, so by Euler’s Formula and the Fundamental Theorem of Arithmetic it suffices to show every positive odd prime is a sum of four squares. [We even know this is true for primes congruent to 1 modulo 4, but will not use this.] Let $p$ be an odd prime. By the lemma and the Well-ordering Principle there exists an element $(x_0, y_0, z_0, w_0, m) \in \mathbb{Z}^5$ with $1 \leq m < p$ and satisfying

\[(*) \quad x_0^2 + y_0^2 + z_0^2 + w_0^2 = mp \]

with $m$ minimal. To finish, it suffices to show that $m = 1$.

Case 1. $m$ is even:

By (*), an even number of the integers $x_0, y_0, z_0, w_0$ must be odd (if any are odd). If precisely two of these are odd, we may assume that they are $x_0, y_0$. So we have $x_0 \equiv y_0 \mod 2$ and $z_0 \equiv w_0 \mod 2$, hence

\[
\frac{(x_0 + y_0)^2}{2} + \frac{(x_0 - y_0)^2}{2} + \frac{(z_0 + w_0)^2}{2} + \frac{(z_0 - w_0)^2}{2} = \frac{m}{2}p,
\]

contradiction the minimality of $m$.

Case 2. $m$ is odd:

We use the modification of the Division Algorithm given by Exercise 4.17(3). Using this modified division algorithm, we have equations:

\[
x_0 = mx_1 + x_2 \quad \text{with } |x_2| < \frac{m}{2},
\]

\[
y_0 = my_1 + y_2 \quad \text{with } |y_2| < \frac{m}{2},
\]

\[
z_0 = mz_1 + z_2 \quad \text{with } |z_2| < \frac{m}{2},
\]

\[
w_0 = mw_1 + w_2 \quad \text{with } |w_2| < \frac{m}{2}.
\]

Therefore, we have

\[
x_2^2 + y_2^2 + z_2^2 + w_2^2 \equiv x_0^2 + y_0^2 + z_0^2 + w_0^2 \equiv 0 \mod m,
\]

so

\[(†) \quad x_2^2 + y_2^2 + z_2^2 + w_2^2 = Mm
\]

for some integer $M$.

Subcase 1. $M = 0$:

In this subcase, we must have $x_2 = y_2 = z_2 = w_2 = 0$ and $m \mid x_0$, $m \mid y_0$, $m \mid z_0$, $m \mid w_0$, so $m^2 \mid x_0^2 + y_0^2 + z_0^2 + w_0^2 = mp$. Hence $m \mid p$ with $1 \leq m < p$. It follows that $m = 1$, so we are done in this subcase.

Subcase 2. $M > 0$:
We have \((x^2 + y^2 + z^2 + w^2) \cdot (x'^2 + y'^2 + z'^2 + w'^2) = A^2 + B^2 + C^2 + D^2\) by Euler’s Equation with \(A = x_1x_2 - y_1y_2 - z_1z_2 - w_1w_2\) and for some integers \(B, C, D\). Consequently,

\[
mp = x'_0^2 + y'_0^2 + z'_0^2 + w'_0^2
= (mx_1 + x_2)^2 + (my_1 + y_2)^2 + (mz_1 + z_2)^2 + (mw_1 + w_2)^2
= m^2(x'^1_1 + y'^1_1 + z'^1_1 + w'^1_1) + 2mA + mM.
\]

Multiply this equation by \(M/m\) and use (†) to obtain the following equation in \(\mathbb{Z}\):

\[
Mp = mM(x^2_1 + y^2_1 + z^2_1 + w^2_1) + 2mA + M^2
= (x^2_2 + y^2_2 + z^2_2 + w^2_2)(x^2_1 + y^2_1 + z^2_1 + w^2_1) + 2mA + M^2
= A^2 + B^2 + C^2 + D^2 + 2mA + M^2 = (A + M)^2 + B^2 + C^2 + D^2
\]

is a sum of four squares with

\[
0 < mM = x^2_2 + y^2_2 + z^2_2 + w^2_2 < 4\left(\frac{m}{2}\right)^2 = m^2.
\]

Therefore, \(0 < M < m\), contradicting the choice of \(m\).

Remark 32.6. It is much harder to determine which positive integers are sums of three squares. This was determined by Gauss, who showed that a positive integer \(n\) is a sum of three squares if and only if \(n \neq 4^m(8k + 7)\) for some non-negative integers \(m\) and \(k\).

[Note that 3 = 1 + 1 + 1 and 13 = 2^2 + 3^2 but 39 = 3 \cdot 13\) is a sum of four squares (as 39 = 1^2 + 2^2 + 3^2 + 5^2), but not fewer.]

Hilbert proved a much more general theorem when he solved the Waring Problem that given any positive integer \(n\) there exists a minimal positive integer \(g(n)\) such that every positive integer is a sum of \(g(n)\) \(n\)th powers. His proof is not constructive, and it is a difficult problem (and open for most \(n\)) to determine \(g(n)\).

Exercises 32.7.

1. Show \(\mathcal{H}\) is a ring and prove quaternion conjugation is multiplicative.
2. Let \(z = x_0 + x_1i + x_2j + x_3k\). We call \(z\) a pure quaternion if \(x_0 = 0\).
   Suppose \(z\) is nonzero. Show that \(z\) is a pure quaternion if and only if \(z\) does not lie in \(\mathbb{Q}\) but \(z^2\) does lie in \(\mathbb{Q}\).
3. Generalize the construction of quaternions as follows: Let \(F\) be any field of characteristic different from two and \(a\) and \(b\) elements of \(F\). Let \(Q\) be the 4-dimensional \(F\)-vector space on basis \(\{1, i, j, k\}\). Make \(Q\) into a ring by defining
   \[
i^2 = a, \quad j^2 = b, \quad ij = -ji = k
   \]
and extend linearly. Define the *conjugate* of $z$ by

$$x_01 + x_1i + x_2j + x_3k := x_01 - x_1i - x_2j - x_3k$$

and the norm map

$$N : Q \to F \text{ by } z \mapsto z \overline{z}$$

and just as for the Hamiltonian quaternions it is multiplicative. Show that

$$N(x_01 + x_1i + x_2j + x_3k) = x_0^2 - ax_1^2 - bx_2^2 + abx_3^2.$$

We call $Q$ a *general quaternion algebra*. In particular, $Q$ is a division ring if it satisfies the property that $N(z) = 0$ if and only if $z = 0$. [If this does not happen, then $Q \cong M_2(F)$.] Can you give non-isomorphic examples, of quaternion algebras that are division rings over $\mathbb{Q}$?

4. Show if an integer $n$ satisfies $n = 4^m(8k + 7)$, then $n$ is not a sum of three squares. [The converse is the hard part.]
CHAPTER VII

Polynomial Rings

In this chapter we study polynomial rings over rings, concentrating on the case of a polynomial ring in one variable over a domain, e.g., a field. Most importantly, we prove Kronecker’s Theorem that a non-constant polynomial over a field $F$ (in one variable) has a root in a field containing $F$. This will be a basic result when we turn to studying Field Theory. We then investigate polynomial rings over a UFD. We prove that any such ring is itself a UFD using a classical result of Gauss. In the final section of this chapter, we attempt to motivate the commutative algebra needed for algebraic geometry that will be studied more systematically in a later chapter, that is the study of zeros of polynomials in many variables over a field, especially algebraically closed fields.

33. Introduction to Polynomial Rings

We start by explicitly defining polynomial rings. If $R$ is a ring, let $R[t] := \{a_0 + a_1t + \cdots + a_nt^n \mid a_i \in R, i = 1, \ldots, n, \text{ some } n \geq 0\}$ with $a_0 + a_1t + \cdots + a_nt^n = 0$ if and only if $a_i = 0$ for $i = 0, \ldots, n$. Therefore, if $f = a_0 + a_1t + \cdots + a_nt^n$ and $g = b_0 + b_1t + \cdots + b_mt^m$, then

$$f = g \text{ if and only if } a_i = b_i \text{ for all } i$$

[where we let $a_i = 0$ for all $i > n$ and $b_j = 0$ for all $j > m$].

With $f$ and $g$ as above, define

$$f + g := (a_0 + b_0) + (a_1 + b_1)t + \cdots + \sum (a_i + b_i)t^i$$

and

$$fg := \sum c_it^i \text{ with } c_i = \sum_{k=0}^{i} a_{i-k}b_k.$$

This makes $R[t]$ into a ring called the polynomial ring over $R$ with $0_{R[t]} = 0_R + 0t + \cdots$ and $1_{R[t]} = 1_R + 0_Rt + \cdots$, where we view $R \subseteq R[t]$ by identifying $r$ and $r1_{R[t]}$ for all $r \in R$, i.e., identifying the ring monomorphism $R \rightarrow R[t]$ given by $r \mapsto r1_{R[t]}$ as the inclusion. So $R$ is
just the set (and subring) of constant polynomials in $R[t]$. Recursively define the polynomial ring in $n$ variables $t_1, \ldots, t_n$ over $R$ by

$$R[t_1, \ldots, t_n] := (R[t_1, \ldots, t_{n-1}])[t_n].$$

If $f = a_0 + a_1 t + \cdots + a_n t^n$ is a nonzero polynomial with $a_n \neq 0$, we define the degree of $f$ by $\deg f := n$ and the leading coefficient of $f$ by lead $f = a_n$. If lead $f = 1$, we say $f$ is monic. If $f = 0$ or $\deg f = 0$, we call $f$ a constant polynomial. We let lead 0 = 0 (and one can define $\deg 0 = -\infty$ if we let $\max(n, -\infty) = n$ and $n + (-\infty) = -\infty$ for any integer $n$).

**Properties 33.1.** Let $R$ be a ring, $f$, $g$ nonzero polynomials.
1. If $R$ is commutative, then so is $R[t]$.
2. $\deg(f + g) \leq \max\{\deg f, \deg g\}$. (Strict inequality is possible, e.g., if $g = -f$.)
3. $\deg(fg) \leq \deg f + \deg g$ with equality if and only if

$$ (*) \quad \text{lead}(fg) = \text{lead}(f) \, \text{lead}(g)$$

is nonzero.
4. If $R$ is a domain, then $(*)$ holds, so $R[t]$ is a domain.

If we had defined the degree of the zero polynomial to be $-\infty$, these properties would still hold for all polynomials.

**Definition 33.2.** Let $S$ be a commutative ring and $a_1, \ldots, a_n$ elements in $S$. If $R$ a subring of $S$, then the map

$$e_{a_1, \ldots, a_n} : R[t_1, \ldots, t_n] \to S \text{ given by } f \mapsto f(a_1, \ldots, a_n)$$

is called evaluation at $a_1, \ldots, a_n$. We denote the image of $e_{a_1, \ldots, a_n}$ by $R[a_1, \ldots, a_n]$, as elements in this image are sums of the form $ra_1^{i_1} \cdots a_n^{i_n}$ with $r \in R$ and $i_1, \ldots, i_n$ non-negative integers. The evaluation map $e_{a_1, \ldots, a_n}$ is a ring homomorphism and if $R = S$, it is surjective. If $R$ is not commutative, then this map may not be a homomorphism. (Why?) For this reason, we shall only study polynomial rings over commutative rings. Of course if the subring of $S$ generated by $R$ and $a_1, \ldots, a_n$ is commutative, the evaluation map $e_{a_1, \ldots, a_n}$ is still a ring homomorphism as we may replace $S$ by a commutative subring.

If $\varphi : R \to S$ is a ring homomorphism (respectively, ring monomorphism, ring epimorphism, ring isomorphism) of commutative rings, then $\varphi$ induces a map

$$\tilde{\varphi} : R[t_1, \ldots, t_n] \to S[t_1, \ldots, t_n]$$

given by $\sum \cdots \sum a_{i_1, \ldots, i_n} t_1^{i_1} \cdots t_n^{i_n} \mapsto \sum \cdots \sum a_{i_1, \ldots, i_n} \varphi(a_{i_1, \ldots, i_n}) t_1^{i_1} \cdots t_n^{i_n}$.
and \( \tilde{\varphi} \) is a ring homomorphism (respectively, ring monomorphism, ring epimorphism, ring isomorphism). For convenience, we write the iterated sum by \( \sum \) or \( \sum_{i_1, \ldots, i_n} \). Moreover, if \( a_1, \ldots, a_n \) are elements in \( R \), we have a commutative diagram:

\[
\begin{array}{ccc}
R[t_1, \ldots, t_n] & \xrightarrow{\varphi} & S[t_1, \ldots, t_n] \\
\downarrow e_{a_1, \ldots, a_n} & & \downarrow e_{\varphi(a_1), \ldots, \varphi(a_n)} \\
R & \xrightarrow{\varphi} & S.
\end{array}
\]

In particular, the above applies to the case that \( S = R/\mathfrak{A} \) with \( \mathfrak{A} \) an ideal in \( R \) and \( \varphi \) the canonical epimorphism \( - : R \to R/\mathfrak{A} \). Assume that \( \mathfrak{A} < R \) (otherwise we would get no further information). Then the canonical epimorphism induces a ring epimorphism, \( \tilde{\varphi} : R[t_1, \ldots, t_n] \to (R/\mathfrak{A})[t_1, \ldots, t_n] \) given by \( \sum_{i_1, \ldots, i_n} a_{i_1, \ldots, i_n} t_1^{i_1} \cdots t_n^{i_n} \mapsto \sum_{i_1, \ldots, i_n} \overline{a_{i_1, \ldots, i_n}} t_1^{i_1} \cdots t_n^{i_n} \). In this case, for simplicity of notation, we write \( - \) for \( \tilde{\varphi} \). Hence if \( (a_1, \ldots, a_n) \in R^n \), we have a commutative diagram:

\[
\begin{array}{ccc}
R[t_1, \ldots, t_n] & \xrightarrow{-} & \overline{R[t_1, \ldots, t_n]} \\
\downarrow e_{a_1, \ldots, a_n} & & \downarrow e_{\overline{a_1}, \ldots, \overline{a_n}} \\
R & \xrightarrow{-} & \overline{R},
\end{array}
\]

i.e.,

\[
\overline{f(x_1, \ldots, x_n)} = f(x_1, \ldots, x_n).
\]

In particular, if no \( (a_1, \ldots, a_n) \in \overline{R}^n \) satisfies \( \overline{f}(a_1, \ldots, a_n) = 0 \), then no \( (x_1, \ldots, x_n) \in R^n \) satisfies \( f(x_1, \ldots, x_n) = 0 \). Viewed geometrically, this says \( f = 0 \) in \( R^n \) has no solution if \( \overline{f} = 0 \) in \( \overline{R}^n \) has no solution. This is a very useful idea, that is often used, for instance in number theory, to try to show a given Diophantine equation has no solution.

We return to the one variable case. For further remarks in this chapter about polynomials in many variables, see Addendum 35.

**Lemma 33.3.** Let \( R \) be a domain. Then \( R[t]^x = R^x \).

**Proof.** Clearly, units in \( R \) remain units in \( R[t] \). Conversely, let \( f, g \) be polynomials in \( R[t] \) satisfying \( fg = 1 \). Then \( \deg(fg) = \deg f + \deg g = 0 \), so \( \deg f = \deg g = 0 \) and \( f, g \) lie in \( R \). \( \square \)

**Note** that \( \mathbb{Z}/4\mathbb{Z} \) is not a domain and \( \overline{1 + 2t}^2 = \overline{1} \) in \( \mathbb{Z}/4\mathbb{Z}[t] \).

By Exercise 29.22(33.4), whose solution mimics how you learned how to divide polynomials, we have the following:
Theorem 33.4. (General Division Algorithm) Let \( f, g \) be polynomials with coefficients in a commutative ring \( R \). Suppose the leading coefficient of \( g \) is a unit. Then there are polynomials \( q \) and \( r \) with coefficients in \( R \) such that \( f = gq + r \) with either \( r = 0 \) or the degree of \( r \) is less than the degree of \( g \).

As usual, if \( R \) is a commutative ring, we say that \( \alpha \) in \( R \) is a root of a polynomial \( f \) in \( R[t] \) if \( f(\alpha) = e_\alpha(f) = 0 \). The General Division Algorithm has many consequences, most that you have encountered in middle school algebra that we now state and prove.

Corollary 33.5. Let \( F \) be a field. Then \( F[t] \) is a euclidean domain. In particular, \( F[t] \) is a PID, hence also a UFD.

Corollary 33.6. (Remainder Theorem) Let \( R \) be a commutative ring, \( \alpha \) an element in \( R \), and \( f \) a nonzero polynomial in \( R[t] \). Then there exists a polynomial \( q \) in \( R[t] \) satisfying \( f = (t-\alpha)q + f(\alpha) \). Moreover,

\[ t-\alpha \mid f \text{ in } R[t] \text{ if and only if } f(\alpha) = 0, \text{ i.e., } \alpha \text{ is a root of } f. \]

Proof. Apply the General Division Algorithm with \( g = t-\alpha \) to get \( f = (t-\alpha)q + r \) in \( R[t] \) for some \( q, r \in R[t] \) with \( r = 0 \) or \( \deg r < \deg(t-\alpha) = 1 \). It follows that \( r \in R \), hence \( f(\alpha) = e_\alpha(f) = e_\alpha((t-\alpha)q) + e_\alpha(r) = r \). The result now follows easily. \( \square \)

Corollary 33.7. Let \( R \) be a domain, \( f \) a nonzero polynomial in \( R[t] \), and \( x_1, \ldots, x_n \) distinct roots of \( f \) in \( R \). Then \( \prod_{i=1}^{n}(t-x_i) \mid f \) in \( R[t] \).

In particular, \( \deg f \geq n \).

Proof. If \( n = 1 \), the result follows by the Remainder Theorem, so assume that \( n > 1 \). By iterating the Remainder Theorem, we can write \( f = (t-x_1)^r h \) in \( R[t] \) for some \( r \in \mathbb{Z}^+ \) and \( h \in R[t] \) with \( h(x_1) \) nonzero. As \( R[t] \) is a domain, we have \( \deg h \leq \deg f - r \). For each \( i = 2, \ldots, n \), applying the evaluation homomorphism \( e_{x_i} \) shows that \( 0 = f(x_i) = (x_i-x_1)^r h(x_i) \) in the domain \( R[t] \). Hence \( h(x_i) = 0 \) for \( i = 2, \ldots, n \). By induction \( \prod_{i=2}^{n}(t-x_i) \mid h \) in \( R[t] \), hence \( (t-x_1)^r \prod_{i=2}^{n}(t-x_i) \mid f \) in \( R[t] \). As \( \deg f \geq \deg h \geq n - 1 \), the result follows. \( \square \)

Note the proof shows that the degree of \( f \) above is greater than the number of roots of \( f \) in \( R \) “counted with multiplicity.”

Corollary 33.8. (Lagrange) Let \( R \) be a domain and \( f \) a nonzero polynomial of degree \( n \) in \( R[t] \). Then \( f \) has at most \( n \) roots in \( R \).

Corollary 33.9. Let \( R \) be a domain and \( f \) and \( g \) nonzero polynomials in \( R[t] \). If \( f(x_i) = g(x_i) \) for \( n \) distinct elements \( x_1, \ldots, x_n \) in \( R \) with \( n > \max(\deg f, \deg g) \), then \( f = g \). In particular, if \( R \) is infinite and \( f(x) = g(x) \) for all \( x \in R \), then \( f = g \).
Remarks 33.10. 1. The polynomial \( t^2 - 1 \) in \((\mathbb{Z}/8\mathbb{Z})[t]\) has four roots: \( \pm 1, \pm 3 \), so the assumption that \( R \) be a domain in the Corollary 33.7 is essential.

2. If \( R \) is not a commutative ring, as mentioned evaluation is not necessarily a ring homomorphism. This occurs because we assume that the variable \( t \) commutes with all elements of \( R \), but the element \( x \) at which we wish to evaluate will, in general, not commute with all elements in \( R \). If \( x \) does commute with all elements of \( R \), then evaluation at that element will be a ring homomorphism. The theory of polynomials over non-commutate rings is, therefore, quite different. For example, let \( \mathcal{H} \) the Hamiltonian quaternions, \( a, b, c \in \mathbb{R} \) satisfy \( 1 = a^3 + b^3 + c^3 \), and \( \alpha = ai + bj + ck \). Then every such \( \alpha \) is a root of the polynomial \( t^2 - 1 \) in \( \mathcal{H}[t] \). So \( t^2 - 1 \) has infinitely many roots in \( \mathcal{H}[t] \).

3. If \( R \) is an infinite domain, Corollary 33.9 says that polynomials in \( R[t] \) and polynomial functions are essentially the same, where \( p : R \to R \) is a polynomial function if there exists a polynomial \( f \) in \( R[t] \) such that \( p(x) = f(x) \) for all \( x \in R \). We write \( p_f \) for this \( p \). So the corollary says that if \( f \) and \( g \) are polynomials in \( R[t] \), with \( R \) an infinite domain, then \( f = g \) if and only if \( p_f = p_g \). (Clearly, if \( f = g \) then \( p_f = p_g \), and the zero polynomial is the only polynomial one with infinitely many roots.)

4. Let \( p \) be a prime in \( \mathbb{Z}^+ \). Then the polynomials \( t^p \) and \( t \) in \((\mathbb{Z}/p\mathbb{Z})[t]\) are clearly distinct, but by Fermat’s Little Theorem, \( p^p = p \) in the notation of the previous remark, so the condition that \( R \) be infinite in the previous remark is crucial.

For the next corollary, we set up important notation. Let \( F \) be a field and \( \mathfrak{A} < F[t] \) be an ideal. We have the canonical ring epimorphism

\[- : F[t] \to F[t]/\mathfrak{A} \text{ given by } g \mapsto \bar{g} = g + \mathfrak{A}.\]

As \( \mathfrak{A} < F[t] \), the ring \( F[t]/\mathfrak{A} \) is not the trivial ring. Since a field is simple, the canonical epimorphism induces a ring monomorphism

\[-|_{F} : F \to F[t]/\mathfrak{A} \text{ given by } a \mapsto \bar{a}.\]

We shall always view this as an inclusion, i.e., for all \( a \in F \), we shall identify \( a \) and \( \bar{a} \).

As \( \bar{t} = t + \mathfrak{A} \), this identification means that the canonical epimorphism maps \( g = \sum a_i t^i \mapsto \bar{g} = \sum a_i \bar{t}^i = \sum a_i \bar{t}^i \), i.e., the canonical epimorphism under this identification is none other than the evaluation map
at \( \overline{t} \), i.e.,
\[(33.11)\quad F[t] \rightarrow F[t]/\mathfrak{A} \text{ is the map } g \mapsto g(\overline{t}).\]

**Corollary 33.12.** (Kronecker’s Theorem) *Let* \( F \) *be a field and* \( f \) *an irreducible element in* \( F[t] \). *Set* \( K = F[t]/(f) \). *Then* \( K \) *is a field. Viewing* \( F \subset K \), *i.e., as a subfield of* \( K \) *as above, we have* \( \overline{t} \) *is a root of* \( f \) *in* \( K \).

**Proof.** As \( f \) is irreducible in the UFD \( F[t] \), it is a prime element. As \( F[t] \) is a PID, the nonzero prime ideal \((f)\) is maximal, so \( K \) is a field. As \( 0 = \overline{f} = f(\overline{t}) \), by the above, the result follows. \( \square \)

**Corollary 33.13.** *Let* \( F \) *be a field and* \( f \) *a non-constant polynomial in* \( F[t] \). *Then there exists a field* \( K \) *with* \( F \subset K \) *a subfield such that* \( f \) *has a root in* \( K \).

**Proof.** We can write \( f = f_1 g \) with \( f_1 \) and \( g \) polynomials in \( F[t] \) and \( f_1 \) irreducible. By Kronecker’s Theorem, \( f_1 \) hence \( f \) has a root in \( F[t]/(f_1) \), a field containing \( F \). \( \square \)

Of course, to make Kronecker’s Theorem and its corollary really useful, we would want to prove results by induction, the natural choice being the degree of a polynomial. That we can do this follows from the following stronger form of Kronecker’s theorem, which we leave as an exercise, but will be essential when we study field theory:

**Proposition 33.14.** *Let* \( F \) *be a field and* \( f \) *a non-constant polynomial of degree* \( n \) *in* \( F[t] \). *Then* \( F[t]/(f) \) *is a vector space over* \( F \) *of dimension* \( n \). *In particular, there exists a field* \( K \) *containing* \( F \) *such that* \( f \) *has a root in* \( K \) *and the dimension of* \( K \) *as a vector space over* \( F \) *is at most* \( n \).

We next prove a fundamental fact about finite fields. To do so, we use the First Sylow Theorem. This result is also important in general field theory.

**Theorem 33.15.** *Let* \( F \) *be a field and* \( G \) *a finite (multiplicative) subgroup of* \( F^\times \). *Then* \( G \) *is cyclic. In particular, if* \( F \) *is a finite field, the group* \( F^\times \) *is cyclic."

**Proof.** Let \( p \) be a prime dividing \( |G| \) and \( P \) a Sylow \( p \)-subgroup of \( G \) (so the only one as \( G \) is abelian). Choose \( x_P \) in \( P \) such that the order of the cyclic subgroup \( \langle x_P \rangle \) generated by \( x_P \) in \( P \) is maximal. Let \( N \) be the order of \( x_P \). Then \( N \) is a power of \( p \) and we must have \( y^N = 1 \) for all \( y \in P \). (Why?) Therefore, the polynomial \( t^N - 1 \) has \( |P| \) distinct roots in \( F \), so \( N = |P| \) by Corollary 33.8. It follows that \( P \) is
cyclic. Let \( x = \prod_{p \mid |G| \text{ a prime}} x_P \in G \). As \( G \) is abelian, it follows from Exercise 13.7(9) that \( G = \langle x \rangle \) is cyclic. \( \square \)

We end this section, with a short discussion of irreducible polynomials over a field \( F \). Clearly, every linear polynomial, i.e., polynomial of degree one in \( F[t] \) is irreducible.

**Definition 33.16.** The following conditions on a field \( F \) are equivalent:

(i) Every non-constant polynomial in \( F[t] \) has a root in \( F \).

(ii) The only irreducible polynomials in \( F[t] \) are linear.

(iii) Every non-constant polynomial in \( F[t] \) factors into a product of linear polynomials.

A field \( F \) satisfying these equivalent conditions is called **algebraically closed**.

When we study field theory, we shall prove (cf 56.12 below):

**Theorem 33.17.** (Fundamental Theorem of Algebra) *The field of complex numbers is algebraically closed.*

We shall, however, assume here that it has already been established.

**Remark 33.18.** The Fundamental Theorem of Algebra allows us to compute all irreducible polynomials over the reals. They are

(i) Linear polynomials in \( \mathbb{R}[t] \).

(ii) *Quadratic polynomials* (i.e., polynomials of degree two) of the form \( at^2 + bt + c \) in \( \mathbb{R}[t] \) satisfying, \( a \neq 0 \) and \( b^2 - 4ac < 0 \).

In particular, any non-constant polynomial in \( \mathbb{R}[t] \) factors into a product of linear and irreducible quadratic polynomials. For example, we have the factorization of the polynomial \( t^4 + 1 = (t^2 + \sqrt{2}t + 1)(t^2 - \sqrt{2}t + 1) \). Over any field \( F \), the polynomial ring \( F[t] \) contains infinitely many non-associative irreducible polynomials. [Can you prove this?]

Finding all irreducible polynomials in \( F[t] \), with \( F \) or \( F[\sqrt{-1}] \) (which is in fact always a field) not algebraically closed is usually very hard, if not impossible, e.g., if \( F = \mathbb{Q} \).

One of our goals will be to show that given any field \( F \), there exists an algebraically closed field \( K \) containing \( F \) and if \( F \neq K \), then no field properly between \( F \) and \( K \) is algebraically closed. Such a field will be called an **algebraic closure** of \( F \). The proof will need Zorn’s Lemma. We shall prove this in Section 50 below. In particular, \( \mathbb{C} \) is not an algebraic closure of \( \mathbb{Q} \).
Exercises 33.19.
1. Let $R$ be a commutative ring. Show that a polynomial $f = a_0 + a_1 t + \cdots + a_n t^n$ in $R[t]$ is a unit in $R[t]$ if and only if $a_0$ is a unit in $R$ and $a_i$ is nilpotent for every $i > 0$.

2. Let $R$ be a nontrivial commutative ring and $f$ a zero divisor in $R[t]$. Show that there exists a nonzero element $b$ in $R$ so that $bf = 0$.

3. Let $R$ be a nontrivial commutative ring. If $f = a_0 + a_1 t + \cdots + a_n t^n$ is a polynomial in $R[t]$ define the formal derivative $f'$ of $f$ to be $f' = a_1 + 2 a_2 t + \cdots + n a_n t^{n-1}$.

(i) Show the usual rules of differentiation hold.

(ii) Suppose $R$ is a field of characteristic zero. Show that a polynomial $f \in R[t]$ is divisible by the square of a non-constant polynomial in $R[t]$ if and only if $f$ and $f'$ are not relatively prime.

4. Let $R$ be a ring and $G$ a group (or monoid). Define

$$R[G] := \{ \sum_{g \in G} a_g g \mid g \in G \text{ and } a_g \in R \text{ almost all } a_g = 0 \}$$

(where almost all zero means that only finitely many are nonzero) with $+$ and $\cdot$ defined by

$$\sum_{g \in G} a_g g + \sum_{g \in G} b_g g = \sum_{g \in G} (a_g + b_g) g \text{ and } \sum_{g \in G} a_g g \cdot \sum_{g \in G} b_g g = \sum_{g \in G} c_g g \text{ where } c_g = \sum_{g = h d} a_h b_d$$

for all $a_g,b_g$ in $G$. Show this is a ring. It is called the group (respectively, monoid) ring of $R$ by $G$. If $\mathbb{N} := \mathbb{Z}^+ \cup \{0\}$, show that $R[\mathbb{N}^n]$ is isomorphic to $R[t_1, \cdots, t_n]$. Describe $R[\mathbb{Z}^n]$.

5. Let $F$ be a subfield of the complex numbers $\mathbb{C}$. Let $f \in F[t]$ be an irreducible polynomial. Show that $f$ has no multiple root in $\mathbb{C}$, i.e., a root $\alpha$ of $f$ satisfying $(t - \alpha)^n \mid f$ in $F[t]$ with $n > 1$.


7. Prove that the conditions in Definition 33.16 are equivalent.

8. Show that the irreducible polynomials in $\mathbb{R}[t]$ are those stated in Remark 33.18.

9. Show that over any field $F$, there exist infinitely many monic irreducible polynomials in $F[t]$. Also show that if $F$ is algebraically closed, then $F$ must have infinitely many elements.

10. Let $F = \mathbb{Z}/p\mathbb{Z}$ with $p$ a prime. Show that $t^4 + 1 \in F[t]$ is reducible.
34. Polynomial Rings over a UFD

In this section, we prove a theorem of Gauss that a polynomial ring over a UFD $R$ is itself a UFD. The idea is to take a polynomial over $R[t]$, with $R$ a UFD, and view it over the UFD $K[t]$ where $K$ is the quotient field of $R$. To make this useful, we must derive a way of pulling information back to $R[t]$.

We have defined the gcd of two nonzero elements in a domain $R$. We can extend this in the obvious way to a finite number of elements, viz., if $a_1, \ldots, a_n$ are elements in $R$ not all zero, then $d$ in $R$ is a *greatest common divisor* or $gcd$ of $a_1, \ldots, a_n$ if it satisfies both of the following:

(i) $d \mid a_i$ for $i = 1, \ldots, n$.
(ii) If $e \mid a_i$ for some $e$ in $R$, $i = 1, \ldots, n$, then $e \mid d$.

[We let the gcd of nonzero $a$ in $R$ and 0 be $a$.]

We have the following lemma that we leave as an exercise.

**Lemma 34.1.** Let $R$ be a UFD and $a_1, \ldots, a_n$ elements in $R$, not all zero. Then a gcd of $a_1, \ldots, a_n$ exists and is unique up to units.

Let $R$ be a UFD and $f = a_n t^n + \cdots + a_0$ a nonzero polynomial in $R[t]$. A gcd of $a_0, \ldots, a_n$ is called a *content* of $f$. It is unique up to units. We call $f$ *primitive* if 1 is a content of $f$. [So a primitive constant polynomial is just a unit in $R$.] For notational convenience, we often use the notation $C(f)$ for a choice of a content of $f$.

**Remarks 34.2.** Let $R$ be a UFD, $b$ a nonzero element of $R$ and $f$ a nonzero polynomial in $R[t]$.

1. If $x$ is an irreducible element in $R$, then it remains irreducible when considered as an element of $R[t]$. (Look at degrees if it would factor.)
2. $C(bf) \approx bC(f)$, where as before $\approx$ means is an associate of.
3. There exists a primitive polynomial $f_1$ in $R[t]$ satisfying $f = C(f) f_1$. (As $R[t]$ is a domain, $\deg f = \deg f_1$.) This follows easily as we can factor out a content.
4. If $\deg f > 0$ and $f$ is irreducible in $R[t]$, then $f$ is primitive. Indeed, we can write $f = C(f) f_1$ for some primitive polynomial $f_1$ in $R[t]$. If $C(f) \neq 1$, then it factors into irreducibles in $R$ hence in $R[t]$. It follows that $f$ cannot be irreducible as $\deg f_1 > 0$, so $f$ is not a unit.
5. $t^2 - 1$ is a primitive polynomial, but not irreducible, so the converse to the previous remark is false.

If $R$ is a UFD, to show $R[t]$ is a UFD, we must show that any nonzero nonunit polynomial in $R[t]$ factors into a product of irreducible polynomials in $R[t]$ and that this factorization is essentially unique.
The existence is more elementary and will follow from the following lemma and its corollary.

**Lemma 34.3.** Let $R$ be a UFD, $K$ the quotient field of $R$, and $f$ a nonzero polynomial in $K[t]$. Then there exist a primitive polynomial $f_1 \in R[t]$ (of the same degree as $f$) and an element $\alpha \in K$ satisfying $f = \alpha f_1$. Moreover, $f_1$ and $\alpha$ are unique up to units in $R$, i.e., up to elements in $R^\times = R[t]^\times$.

**Proof.** **Existence:** Write $f = \sum_{i=0}^{n} \frac{a_i}{b_i} t^i$ with $a_i, b_i \in R$ and $b_i \neq 0$, $i = 0, \ldots, n$. We clear denominators. Let $b = b_0 \cdots b_n \neq 0$ in the domain $R$. Then $bf \in R[t]$. Let $c = C(bf)$. Then there exists a primitive polynomial $f_1$ in $R[t]$ such that $bf = cf_1$, so $f = \frac{c}{b} f_1$ in $K[t]$ as needed.

**Uniqueness** Suppose that $\frac{c}{b} f_1 = f = \frac{d}{e} f_2$ in $K[t]$, with nonzero $c, d, e \in R$ and $f_1, f_2$ primitive in $R[t]$. Then in $R[t]$, we have $c e f_1 = b d f_2$. Therefore, $c e \approx C(c e f_1) \approx C(b d f_2) \approx b d$ in $R$. Consequently, there exists a unit $u$ in $R^\times$ satisfying $c e u = b d$ in $R$, hence $\frac{c}{b} = u \frac{d}{e}$ in $K$. Thus $\frac{d}{e} f_2 = \frac{c}{b} f_1 = u \frac{d}{e} f_1$ in the domain $K[t]$. It follows that $f_2 = u f_1$. \hfill \Box

**Corollary 34.4.** Let $R$ be a UFD, $f, g$ primitive polynomials in $R[t]$, $h$ a polynomial in $R[t]$, and $K$ the quotient field of $R$.

1. If $g = sf$ in $K[t]$ for some $s$ in $K$, then $s$ is a unit in $R$.
2. If $g = ah$ in $R[t]$ for some $a$ in $R$, then $a$ is a unit in $R$ and $h$ is primitive in $R[t]$.

In particular, if $g$ is non-constant primitive polynomial in $R[t]$ but not irreducible, then $g = q_1 q_2$ for some $q_1, q_2$ in $R[t]$ satisfying $0 < \deg q_i < \deg g$, for $i = 1, 2$.

**Proof.** (1): We have $1 \cdot g = s \cdot f$ in $K[t]$ with both $f$ and $g$ primitive, so $s = u \cdot 1$ for some unit $u$ in $R$ by Lemma 34.3, hence $s$ is a unit in $R$.

(2): Write $h = C(h) h_1$ with $h_1$ primitive in $R[t]$. Then $1 \cdot g = a \cdot C(h) h_1$, so $a C(h) \approx 1$, hence $a$ is a unit in $R$ and $h$ is primitive. [Note if $x$ is a nonzero element in $R$ then $x/x$ is a unit but $x$ is not necessarily a unit.] \hfill \Box

The key to showing the uniqueness statement for factorization in $R[t]$ when $R$ is a UFD is the following:
Lemma 34.5. (Gauss’ Lemma) Let $R$ be a UFD and $f, g$ non-constant polynomials. Then $C(fg) \approx C(f)C(g)$. In particular, the product of primitive polynomials in $R[t]$ is primitive.

Proof. Write $f = C(f)f_1$ and $g = C(g)g_1$ with $f_1$ and $g_1$ primitive polynomials in $R[t]$. Then

$$C(fg) \approx C(f)f_1C(g)g_1 \approx C(f)C(g)f_1g_1$$

and

$$C(f)C(g) \approx C(f)f_1C(g)g_1 \approx C(f)C(g)f_1C(g_1),$$

so it suffices to show the last statement, i.e., we may assume that $f$ and $g$ are primitive in $R[t]$ and must show $fg$ is primitive in $R[t]$. Suppose this is false. Then there exists an irreducible element $p$ such that $p \mid C(fg)$ in $R$. As $R$ is a UFD, $p$ is a prime element, hence $\overline{R} = R/(p)$ is a domain. Let $\overline{f} : R[t] \to (R/(p))[t]$ be the ring epimorphism given by $\sum a_i t^i \mapsto \sum \overline{a_i} \overline{t^i}$. By assumption, $\overline{0} = \overline{f}g = \overline{f}g$ in the domain $\overline{R}[t]$. As $f$ and $g$ are primitive in $R[t]$, the prime $p$ does not divide some coefficient of $f$ and some coefficient of $g$, i.e., $\overline{f} \neq 0$ and $\overline{g} \neq 0$ in $\overline{R}[t]$, contradicting the fact that $\overline{R}[t]$ is a domain.

Corollary 34.6. Suppose that $R$ is a UFD with quotient field $K$. Let $f, g$ be non-constant primitive polynomials in $R[t]$ and $h$ a non-constant polynomial in $K[t]$. If $f = gh$ in $K[t]$, then $h$ lies in $R[t]$ and is primitive.

Proof. By Lemma 34.3, we can write $h = \alpha h_1$ with $\alpha \in K$ and $h_1$ a primitive polynomial in $R[t]$, so $f = \alpha gh_1$. By Gauss’ Lemma, $gh_1$ is primitive, so $\alpha$ lies in $R^\times$ by Corollary 34.4. It follows that $h$ lies in $R[t]$ and is primitive.

A similar proof yields:

Lemma 34.7. Let $R$ be a UFD with quotient field $K$ and $f$ a non-constant irreducible polynomial in $R[t]$. Then $f$ remains irreducible in $K[t]$.

Proof. As $f$ is a non-constant irreducible polynomial in $R[t]$, it is primitive. Suppose that $f$ factors as $f = g_1g_2$ in $K[t]$, with $g_i \in K[t]$ and $0 < \deg g_i < \deg f$, for $i = 1, 2$. By Lemma 34.3, we can write $g_i = \alpha_i h_i$ with $\alpha_i \in K$ and $h_i$ in $R[t]$ primitive (and $\deg g_i = \deg h_i$), for $i = 1, 2$. Therefore, we have $f = (\alpha_1\alpha_2)h_1h_2$. By Gauss’ Lemma, $h_1h_2$ is primitive, so by Corollary 34.4, we have $\alpha = \alpha_1\alpha_2$ is a unit in $R$ and $f = \alpha h_1h_2$ in $R[t]$. It follows that $f$ is reducible in $R[t]$, a contradiction.
We can now prove our theorem.

**Theorem 34.8.** Let \( R \) be a UFD. Then \( R[t] \) is a UFD.

**Proof.** Existence: Every nonzero element in \( R \) is a unit or factors into irreducibles that remain irreducible in \( R[t] \), so it suffices to factor a non-constant polynomial in \( R[t] \). Let \( f \) be such a polynomial. As \( f = C(f)f_1 \) with \( f_1 \) a primitive polynomial in \( R[t] \), we may also assume that \( f \) is primitive. If \( f \) is not irreducible, then \( f = g_1g_2 \) for some polynomials \( g_i \in R[t] \) satisfying \( 0 < \deg g_i < \deg f \) for \( i = 1, 2 \) by Corollary 34.4. By induction on \( \deg f \), we conclude that \( g_1 \) and \( g_2 \) hence \( f \) factor into irreducible polynomials in \( R[t] \).

Uniqueness: We know that any nonzero element in \( R \) has a unique factorization (up to associates and order), so the content of a polynomial factors uniquely up to units. It follows by Lemma 34.3 that it suffices to show any non-constant primitive polynomial \( f \) has a unique factorization (up to associates and order). Suppose that \( f_1 \cdots f_r = f = g_1 \cdots g_s \) with all the \( f_i \) and \( g_j \) irreducible in \( R[t] \). By Lemma 34.7, all the \( f_i \) and \( g_j \) remain irreducible in \( K[t] \), where \( K \) is the quotient field of \( R \). As \( K[t] \) is a UFD, \( r = s \), and after relabeling, we may assume that \( f_i = c_ig_i \) in \( K[t] \) for some \( c_i \in K^\times \) for each \( i \). But the \( f_i \) and \( g_j \) are primitive in \( R[t] \) as they are irreducible, so \( c_i \in R^\times \) for all \( i \) by Corollary 34.4. We conclude that \( f_i \approx g_i \) in \( R[t] \) for all \( i \), and the proof is complete. \( \square \)

**Corollary 34.9.** If \( R \) is a UFD, then so is \( R[t_1, \ldots, t_n] \).

**Examples 34.10.** Let \( R \) be a UFD and \( F \) a field.

1. \( \mathbb{Z}[t] \) and \( F[t_1, t_2] \) are UFDs but not PIDs. In fact, \( R[t] \) is a PID if and only if \( R \) is a field.

2. A polynomial ring in any number of variables [even infinitely many — definition of such a polynomial ring?] over \( R \) is a UFD.

3. Let \( \varphi : R \to S \) be a ring epimorphism of domains. Then \( S \) is not necessarily a UFD (even if a domain), e.g. \( e_{\sqrt{-5}} : \mathbb{Z}[t] \to \mathbb{Z}[\sqrt{-5}] \) is such an example.

We end this discussion with a few remarks about irreducible elements in \( R[t] \), with \( R \) a domain.

**Remarks 34.11.** 1. (Eisenstein’s Criterion) Let \( R \) be a UFD and \( K \) its quotient field. Let \( f = \sum_{i=0}^{n} a_it^i \) be a non-constant polynomial in \( R[t] \) of degree \( n \). Suppose that there exists an irreducible element \( p \) in \( R \) (hence a prime) satisfying:
4. Let $R$.

2. Let $R$.

1. Prove Lemma 34.1.

Exercises 34.12.

1. Prove Lemma 34.1.

2. Let $R$ be a domain that is not a field. Show that $R[t]$ is not a PID.

3. Let $R$ be a UFD, $K$ its quotient field. Let $f$ and $g$ be non-constant polynomials in $R[t]$. Write $f = C(f) f_1$ and $g = C(g) g_1$ with $f_1$ and $g_1$ primitive polynomials in $R[t]$. Show

   (i) If $f|g$ in $K[t]$ then $f_1|g_1$ in $R[t]$. In particular, if $f$ and $g$ are primitive, then $f|g$ in $K[t]$ if and only if $f|g$ in $R[t]$.

   (ii) Suppose that $f$ and $g$ are primitive. Then $f$ and $g$ have a common factor over $K[t]$ if and only if they have a common factor over $R[t]$.

4. Let $R$ be a commutative ring, and $\{t_i, i \in I\}$ indeterminants. Define the polynomial ring $R[t_i]_{i \in I}$ in the variables $\{t_i, i \in I\}$. Show that it is a UFD if $R$ is.
5. Let \( f = \sum_{i=0}^{n} a_i t^i \) be a polynomial in \( \mathbb{Z}[t] \) with \( a_n \neq 0 \). Let \( r = a/b \) with \( b \) nonzero and \( a \) and \( b \) relatively prime integers. If \( r \) is a root of \( f \), show that \( b \mid a_n \) and if \( a \) is also nonzero, then \( a \mid a_0 \). In particular, if \( f \) is monic, any rational root of \( f \), if any, is an integer.

6. Prove Eisenstein’s Criterion (Remark 34.11(1)).

7. Let \( y = t + a \). Show that the map \( \mathbb{R}[t] \to \mathbb{R}[y] \) given by \( \sum a_i t^i \mapsto \sum a_i y^i \) is a ring isomorphism.

8. Prove the identity of binomial coefficients

\[
\binom{n+1}{n-j} = \sum_{i=j}^{n} \binom{i}{i-j}
\]

and use it to prove that the polynomial

\[
t^{p^r-1}(p-1) + t^{p^r-1}(p-2) + \cdots + t^{p^r-1} + 1
\]

is an irreducible polynomial in \( \mathbb{Z}[t] \) (hence \( \mathbb{Q}[t] \)) for every (positive) prime \( p \) in \( \mathbb{Z} \).

35. Addendum: Polynomial Rings over a Field

In this addendum, we make some motivational remarks that we shall come back to when studying modules. If \( R \) is a commutative ring, we have defined the ring \( R[t_1, \ldots, t_n] \). We view \( R \) as a subring of \( R[t_1, \ldots, t_n] \) of constant polynomials. We look at the case, when \( R \) is a field.

Let \( F \) be a field and \( S \) be a nontrivial commutative ring. We have seen that any ring homomorphism \( \varphi : F \to S \) is monic. Let \( \mathfrak{A} < F[t_1, \ldots, t_n] \) be an ideal. We apply this to \( S = F[t_1, \ldots, t_n]/\mathfrak{A} \). Let \( \overline{\cdot} : F[t_1, \ldots, t_n] \to F[t_1, \ldots, t_n]/\mathfrak{A} \) be the canonical epimorphism. Then the composition

\[
(*) \quad F \subset F[t_1, \ldots, t_n] \xrightarrow{\overline{\cdot}} F[t_1, \ldots, t_n]/\mathfrak{A}
\]

is also monic. As before, we view this as an embedding, i.e., identify \( a \) in \( F \) with \( \overline{a} \) in \( F[t_1, \ldots, t_n]/\mathfrak{A} \). Under this identification the map \( \overline{\cdot} : F[t_1, \ldots, t_n] \to F[t_1, \ldots, t_n]/\mathfrak{A} \) becomes \( \overline{f} = f(\overline{t_1}, \ldots, \overline{t_n}) \), so it is the evaluation map \( e_{(\overline{t_1}, \ldots, \overline{t_n})} : F[t_1, \ldots, t_n] \to F[t_1, \ldots, t_n]/\mathfrak{A} \).

Observe that if \( \mathfrak{A} \subset \mathfrak{B} < F[t_1, \ldots, t_n] \) are ideals, then we have a
commutative diagram

\[
\begin{array}{ccc}
F[t_1, \ldots, t_n] & \longrightarrow & F[t_1, \ldots, t_n]/\mathfrak{B} \\
\downarrow & & \downarrow \\
F[t_1, \ldots, t_n]/\mathfrak{A} & \end{array}
\]

with all the maps the obvious ring epimorphisms. The most interesting case is when \(\mathfrak{B}\) is a maximal ideal, hence \(F[t_1, \ldots, t_n]/\mathfrak{B}\) is a field. We look at a special case of this.

Let \(x = (x_1, \ldots, x_n)\). For each \(x = (x_1, \ldots, x_n)\) in \(F^n\) set

\[m_x = (t_1 - x_1, \ldots, t_n - x_n),\]

an ideal in \(F[t_1, \ldots, t_n]\). Let \(\mathfrak{A} = m_x\) in (*) . For each \(i = 1, \ldots, n\), we have \(t_i - x_i\) lies in \(F\), so we can view \(F[t_1, \ldots, t_n]/m_x \subset F\). By (*), we conclude that \(F[t_1, \ldots, t_n]/m_x = F\). In particular, \(m_x\) is a maximal ideal in \(F[t_1, \ldots, t_n]\). We conclude that the canonical epimorphism

\[- : F[t_1, \ldots, t_n] \rightarrow F[t_1, \ldots, t_n]/m_x\]

takes \(f \in F[t_1, \ldots, t_n]\) to \(\overline{f} = f(x) = e_x(f)\) in \(F\). So we have

\[f(x) = 0 \text{ in } F \text{ if and only if } \overline{f} = 0 \text{ in } F \text{ if and only if } f \in m_x.\]

The question of interest is to solve the following:

**Problem 35.1.** Let \(f_1, \ldots, f_r\) be polynomials in \(F[t_1, \ldots, t_n]\). Does there exist an \(x\) in \(F^n\) such that \(f_i(x) = 0\) for \(i = 1, \ldots, r\), i.e., a common point lying on all the hypersurfaces \(f_1 = 0, \ldots, f_r = 0\) in \(F^n\)?

Let \(\mathfrak{A} = (f_1, \ldots, f_r)\), the ideal generated by the \(f_i\) in \(F[t_1, \ldots, t_n]\) and \(x \in F^n\). Then we have \(f_i(x) = 0\) for \(i = 1, \ldots, r\) if and only if \(f(x) = 0\) for all \(f \in \mathfrak{A}\). We set up a useful notation. If \(\mathfrak{B}\) is an ideal in \(F[t_1, \ldots, t_n]\), let

\[Z_F(\mathfrak{B}) := \{x \in F^n \mid f(x) = 0 \text{ for all } f \in \mathfrak{B}\}.\]

We call \(Z_F(\mathfrak{B})\) the **affine variety** defined by \(\mathfrak{B}\) in \(F^n\). If \(\mathfrak{B} = (g_1, \ldots, g_s)\), we write \(Z_F(g_1, \ldots, g_s)\) for \(Z_F(\mathfrak{B})\).

So our problem is: Is \(Z_F(\mathfrak{A})\) non-empty?

If \(\mathfrak{A}\) is the unit ideal, then \(Z_F(\mathfrak{A})\) is empty as \(f(x) = 1\) if \(f\) is the constant function 1. So a necessary condition is that \(\mathfrak{A} < F[t_1, \ldots, t_n]\). The equation \(t_1^2 + t_2^2 = -1\) has no solution in \(\mathbb{R}^2\), so \(Z_F(t_1^2 + t_2^2 + 1)\) is empty. In particular, \((t_1^2 + t_2^2 + 1)\) lies in no \(m_x\) with \(x\) in \(\mathbb{R}^n\). Therefore, in general, the answer is no, and we would need some stronger condition on the field \(F\).
Next suppose that \( \mathfrak{A} \) is an ideal in \( F[t_1, \ldots, t_n] \) such that there exists an \( \bar{x} \in F^n \) with \( \mathfrak{A} \subset m_{\bar{x}} \). In particular, if \( f \in \mathfrak{A} \), then \( f \in m_{\bar{x}} \), i.e., if \( \mathfrak{A} \subset m_{\bar{x}} \), then \( \bar{x} \in Z_F(\mathfrak{A}) \). Conversely, suppose that \( \bar{x} \in Z_F(\mathfrak{A}) \). If \( f \in \mathfrak{A} \), then we have \( f(\bar{x}) = 0 \), hence \( f \) lies in the kernel of the map \( e_{\bar{x}} : F[t_1, \ldots, t_n] \to F[t_1, \ldots, t_n]/m_{\bar{x}} \), i.e., \( f \in m_{\bar{x}} \). Therefore,

\[
(35.2) \quad \mathfrak{A} \subset m_{\bar{x}} \quad \text{if and only if} \quad \bar{x} \in Z_F(\mathfrak{A}).
\]

This says that the answer to our problem will be yes for every ideal \( \mathfrak{A} \subset F[t_1, \ldots, t_n] \), if for each ideal \( \mathfrak{A} \), there exists an \( \bar{x} \in F^n \) such that \( \mathfrak{A} \subset m_{\bar{x}} \). In particular, this will be true if and only if whenever \( m \) is a maximal ideal in \( F[t_1, \ldots, t_n] \), there exists an \( \bar{x} \in F^n \) satisfying \( m = m_{\bar{x}} \). (Cf. Exercise 25.21(13).) [Recall we have shown (using Zorn’s Lemma) that every nonunit ideal in a nontrivial commutative ring lies in a maximal ideal (for a Noetherian ring this would follow by the Maximal Principle).]

Hilbert proved the following:

**Theorem 35.3.** (Hilbert Basis Theorem) Let \( R \) be a Noetherian ring. Then \( R[t_1, \ldots, t_n] \) is Noetherian. In particular, if \( F \) is a field then \( F[t_1, \ldots, t_n] \) is Noetherian.

This means that if \( \mathfrak{A} \) is ideal in \( F[t_1, \ldots, t_n] \) with \( F \) a field, then \( Z_F(\mathfrak{A}) = Z_F(f_1, \ldots, f_r) \), for some \( f_1, \ldots, f_r \) in \( F[t_1, \ldots, t_r] \). He also proved the following wonderful theorem:

**Theorem 35.4.** (Hilbert Nullstellensatz) Let \( F \) be an algebraically closed field. Then every maximal ideal in \( F[t_1, \ldots, t_n] \) is of the form \( m_{\bar{x}} \) for some \( \bar{x} \) in \( F^n \). In particular, if \( \mathfrak{A} < F[t_1, \ldots, t_n] \) is an ideal, then \( Z_F(\mathfrak{A}) \) is not empty.

The above is also called the Weak Hilbert Nullstellensatz. It says, under the additional requirement that \( F \) be algebraically closed, the map

\[
F^n \to \{ m \mid m < F[t_1, \ldots, t_n] \text{ a maximal ideal} \} \quad \text{given by} \quad \bar{x} \mapsto m_{\bar{x}}
\]

is a bijection. Actually, Hilbert proved more. In the above, let

\[
\sqrt{\mathfrak{A}} = \{ f \in F[t_1, \ldots, t_n] \mid f^n \in \mathfrak{A} \text{ for some } n \in \mathbb{Z}^+ \},
\]

the radical of the ideal \( \mathfrak{A} \). Then the Hilbert Strong Nullstellensatz says, if \( F \) is algebraically closed, \( \mathfrak{A} < F[t_1, \ldots, t_n] \) an ideal, then

\[
\sqrt{\mathfrak{A}} = \bigcap_{\mathfrak{A} \subset m_{\bar{x}}} m_{\bar{x}}.
\]

We will return to this later and prove these assertions.
Exercise 35.5. Let $R = F[t_1, \ldots, t_n]$, $F$ a field, and $\mathfrak{A}$, $\mathfrak{B}$, and $fA_i$, $i \in I$, be ideals in $R$. Show the following:

(i) If $\mathfrak{A} \subset \mathfrak{B}$, then $Z_F(\mathfrak{B}) \subset Z_F(\mathfrak{A})$.
(ii) $Z_F(\emptyset) = F^n$.
(iii) $Z_F(R) = \emptyset$.
(iv) $Z_F(\sum \mathfrak{A}_i) = \bigcap I(\mathfrak{A}_i)$.
(v) $Z_F(\mathfrak{A}\mathfrak{B}) = Z_F(\mathfrak{A} \cap \mathfrak{B}) = Z_F(\mathfrak{A}) \cup Z_F(\mathfrak{B})$.
(vi) $Z_F(\mathfrak{A}) = Z_F(\sqrt{\mathfrak{A}})$.

36. Addendum: Algebraic Weierstraß Preparation Theorem

If $R$ is a commutative ring, we have defined the ring of formal power series $R[[t]]$. Inductively, let $R[[t_1, \ldots, t_n]] := (R[[t_1, \ldots, t_{n-1}]])(t_n)$, the formal power series in the variables $t_1, \ldots, t_n$.

Previous exercises imply that $R[[t_1, \ldots, t_n]] \times R = R + (t_1, \ldots, t_n)$, i.e., formal power series with unit constant term, (cf. Exercise 25.21(2)) and if $R$ is a domain, so is $R[[t_1, \ldots, t_n]]$ (cf. Exercise 25.21(1)) Moreover, if $R$ is Noetherian, so is $R[[t_1, \ldots, t_n]]$ by Example 29.13(2).

Further, as with polynomials, if $a_1, \ldots, a_n$ lie in $R$, then the evaluation map $e_{a_1, \ldots, a_n} : R[[t_1, \ldots, t_n]] \to R$ by $f \mapsto f(a_1, \ldots, a_n)$ is a ring epimorphism.

In this section we shall show that if $F$ is a field then $F[[t_1, \ldots, t_n]]$ is a UFD. By Exercise 36.11(2), it is also a local ring, i.e., commutative ring with a unique maximal ideal. One nice thing about formal power series is that we do not have to worry about convergence. This means that proving the algebraic analogue of the Weierstraß Preparation Theorem in complex analysis becomes a formal proof, hence much easier.

Definition 36.1. Let $F$ be a field and $f$ an element of $F[[t_1, \ldots, t_n]]$. We say $f$ is regular in $t_n$ if $f(0, \ldots, 0, t_n) \neq 0$. So if $f$ is regular in $t_n$, the element $f$ has a term $ct_n^s$ for some nonzero element $c$ of $F$. If $f$ is regular, write

$$f = g_0 + g_1 t_n + g_2 t_n^2 + \cdots$$

with $g_i \in F[[t_1, \ldots, t_{n-1}]]$ for all $i$ and satisfying

$$0 \neq f(0, \ldots, 0, t_n) = g_0(0, \ldots, 0) + g_1(0, \ldots, 0) t_n + g_2(0, \ldots, 0) t_n^2 + \cdots$$

If $s$ is the first $i$ such that $g_i(0, \ldots, 0)$ is nonzero, we say $s$ is the order of $t_n$ in $f(0, \ldots, 0, t_n)$ and write $s = \text{ord}_{t_n} f$.

We need two lemmas.
Lemma 36.2. Let \( g \) be a nonzero element in \( F[[t_1, t_2]] \). Then there exists a positive integer \( r \) satisfying \( g(t_2^r, t_2) \) is nonzero.

**Proof.** Order the nonzero monomials \( a_{r_1 r_2} t_1^{r_1} t_2^{r_2} \) of \( g \) lexicographically, i.e., \((r_1, r_2) \geq (s_1, s_2)\) if \( r_1 > s_1 \) or \( r_1 = s_1 \) and \( r_2 \geq s_2 \). Let \( a_{s_1 s_2} \) be the smallest monomial in \( g \) relative to this ordering. Choose \( r > s_2 \) and let \( t_1 \mapsto t_2^r \). If \((r_1, r) > (s_1, s_2)\), then \( r r_1 > r s_1 + s_2 \) since either \( r_1 > s_1 \) (and \( r > s_2 \)) or \( r_1 = s_1 \) and \( r_2 > s_2 \). The only term of order \( r s_1 + s_2 \) in \( g(t_2^r, t_2) \) is \( a_{s_1 s_2} t_2^{r s_1 + s_2} \). \( \square \)

The next lemma shows that the notion of regularity is rather mild.

Lemma 36.3. Let \( f_1, \ldots, f_m \) be nonzero elements in \( F[[t_1, \ldots, t_n]] \). Then there exists a ring automorphism \( \sigma \) of \( F[[t_1, \ldots, t_n]] \) fixing \( F[[t_n]] \) with \( \sigma(f_i) \) regular in \( t_n \) for \( i = 1, \ldots, m \).

**Proof.** As \( F[[t_1, \ldots, t_n]] \) is a domain, \( f = f_1 \cdots f_m \) is nonzero. So we are done if we find a \( \sigma \) that works for \( f \), i.e., we may assume that \( f = f_1 \). By the previous lemma, there exist \( r_1, \ldots, r_{n-1} \) such that \( f(t_1^{r_1}, t_2, \ldots, t_n) \neq 0 \), \( f(t_1^{r_2}, t_2^{r_2}, t_3, \ldots, t_n) \neq 0 \), \ldots, \( f(t_1^{r_{n-1}}, t_2^{r_{n-1}}, \ldots, t_n) \neq 0 \). Let \( \sigma \) be the ring automorphism of \( F[[t_1, \ldots, t_n]] \) fixing \( F[[t_n]] \) induced by \( t_1 \mapsto t_1 + t_n^{r_i} \) for \( i = 1, \ldots, n-1 \). [Why does it exist? Its inverse is \( t_1 \mapsto t_1 - t_n^{r_i} \).] Then \( \sigma(f(t_1, \ldots, t_n)) = f(t_1 + t_n^{r_1}, \ldots, t_{n-1} + t_1^{r_{n-1}}, t_n) \). Consequently, \( c_{(0, \ldots, 0, t_n)}(\sigma(f)) = f(t_1^{r_1}, t_2^{r_2}, \ldots, t_1^{r_{n-1}}, t_n) \neq 0 \). \( \square \)

Theorem 36.4. (Algebraic Weierstraß Preparation Theorem) Let \( R = F[[t_1, \ldots, t_n]] \) with \( F \) a field and \( f \in R \) regular in \( t_n \). Suppose that \( s = \text{ord}_{t_n} f \) and \( g \) lies in \( R \). Then there exist unique elements \( h \) and \( r \) in \( R \) satisfying \( g = h f + r \) with \( r \in F[[t_1, \ldots, t_{n-1}]] [t_n] \) and \( \text{deg}_{t_n} r < s \).

**Proof.** We induct on \( n \).

Suppose that \( n = 1 \) and \( t = t_1 \): Then

\[
 f = c_s t^s + \sum_{j > s} c_j t^j \quad \text{with} \quad c_s \neq 0
\]

\[
 g = \sum_{i=0}^{s-1} a_i t^i + \sum_{i \geq s} a_i t^i = \sum_{i=0}^{s-1} a_i t^i + t^s \left( \sum_{i \geq s} a_i t^{i-s} \right).
\]

Let

\[
 r = \sum_{i=0}^{s-1} a_i t^i \quad \text{and} \quad h = \left( \sum_{i \geq s} a_i t^{i-s} \right) \left( \sum_{j \geq s} c_j t^{j-s} \right)^{-1}.
\]
[Note that \(c_s \neq 0\), so the second power series in \(h\) is a unit by Exercise 25.21(2).] Then \(r, h\) satisfy \(g = hf + r\). Since \(hf = 0\) or \(\text{ord}_t(hf) \geq s\), we see that \(r = \sum_{i=0}^{s-1} a_i t^i\) is uniquely determined and hence so is \(h\).

Suppose that \(n > 1\): The result holds for \(R_0 = F[[t_2, \ldots, t_n]]\) by induction. Write

\[
f = \sum_{i=0}^{\infty} f_i t_1^i \quad \text{and} \quad g = \sum_{i=0}^{\infty} g_i t_1^i
\]

with \(f_i, g_i\) in \(R_0\) for all \(i\). Set \(h = \sum_{i=0}^{\infty} h_i t_1^i\) and \(r = \sum_{i=0}^{\infty} r_i t_1^i\) with all the \(h_i\) and \(r_i\) in \(R_0\) to be determined by the equation

\[
\sum_{i=0}^{\infty} r_i t_1^i = \sum_{i=0}^{\infty} g_i t_1^i - (\sum_{i=0}^{\infty} h_i t_1^i) (\sum_{i=0}^{\infty} f_i t_1^i),
\]

i.e.,

\[
\begin{align*}
  r_0 &= g_0 - h_0 f_0 \\
  r_1 &= g_1 - (h_0 f_1 + h_1 f_0) \\
  &\vdots \\
  r_i &= g_i - (h_0 f_i + \cdots + h_i g_0) \\
  &\vdots
\end{align*}
\]

Since \(f(0, \ldots, 0, t_n) = f_0(0, \ldots, t_n)\) and \(f_0\) in \(R_0\) is regular in \(t_n\) with \(s = \text{ord}_{t_n} f\), by induction there exist unique \(h_0, r_0\) with the appropriate properties, hence unique \(r_1, h_1, \ldots, \) etc. working from the top of (*) down with \(\text{deg}_{t_n} r_i < s\) for each \(i\), where we have replaced \(f\) by \(f_0\) and \(g\) by \(g_0 - h_0 f_1\), etc. As all the \(h_i\), \(r_i\) are unique, so are \(r\) and \(h\). \(\square\)

We wish to prove that \(F[[t_1, \ldots, t_n]]\) is a UFD when \(F\) is a field by reducing to the polynomial case where we can use the fact that a polynomial ring over a UFD is a UFD. We begin with the following definition.

**Definition 36.5.** Let \(f \in F[[t_1, \ldots, t_n]]\) with \(F\) a field. Then \(f\) is called a pseudo-polynomial of degree \(s\), if \(f = t_n^s + a_1 t_n^{s-1} + \cdots + a_s\) with \(a_i\) in \(F[[t_1, \ldots, t_{n-1}]]\) and satisfying \(a_i(0, \ldots, 0) = 0\) for \(1 \leq i \leq s\).

In this language, the Weierstraß Preparation Theorem implies:

**Corollary 36.6.** Let \(R = F[[t_1, \ldots, t_n]]\) with \(F\) a field and \(f\) an element of \(R\) regular in \(t_n\) with \(s = \text{ord}_{t_n} f\). Then there exists a unique pseudo-polynomial of degree \(s\) that is an associate of \(f\) in \(F[[t_1, \ldots, t_n]]\).
Theorem 36.9. Let \( R = F[[t_1, \ldots, t_{n-1}]] \) with \( \deg_{t_n} r < s \) satisfying \( t_n^s = hf - r \). Write \( r = \sum_{i=0}^{s-1} r_i t_i^s \) with \( r_i \) in \( F[[t_1, \ldots, t_{n-1}]] \) for all \( i \).

Claim. \( f^* = t_n^s + r \) works.

Since \( f^* = hf \) with \( h \) and \( r \) unique, it suffices to show that \( r_i(0, \ldots, 0) = 0 \) for all \( i \) and \( h \in R^\times \). For the latter, by Exercise 25.21(2), it suffices to show that \( h(0, \ldots, 0) \) is nonzero. But if it is, then \( hf \) has no monomial of the form \( at_n^s \) with \( a \in F^\times \), contradicting \( f^* = t_n^s + r \). Therefore, \( h \) is a unit. Write \( f = c_s t_n^s + \sum_{j>s} f_j t_n^j \) with \( c_s \) nonzero in \( F \) and each \( f_j \in F[[t_1, \ldots, t_{n-1}]] \). Then

\[
h(0, \ldots, 0, t_n) \left( c_s t_n^s + \sum_{j>s} f_j (0, \ldots, 0) t_n^j \right) = t_n^s + \sum_{i=0}^{s-1} r_i (0, \ldots, 0) t_n^i.
\]

Comparing \( t_n^s \) terms shows that \( r_i(0, \ldots, 0) = 0 \) for all \( i \).

Notation 36.7. We denote the unique pseudo-polynomial of \( f \) in Corollary 36.6 by \( f^* \).

Corollary 36.8. Let \( R = F[[t_1, \ldots, t_n]] \) with \( F \) a field and \( f \) and \( g \) elements of \( R \) both regular in \( t_n \). Then \((fg)^* = f^*g^*\).

Proof. Note that \( fg \) is regular at \( t_n \) as \( f \) and \( g \) are. By the previous result, there exist unique units \( u, v, w \in R^\times \) satisfying \( f^* = uf, g^* = vg \), and \((fg)^* = wfg \) pseudo-polynomials of the appropriate degree. Thus \( f^*g^* = uvfg \) and uniqueness shows that \((fg)^* = f^*g^*\).

Theorem 36.9. Let \( F \) be a field. Then \( F[[t_1, \ldots, t_n]] \) is a UFD.

Proof. We induct on \( n \), the case \( n = 0 \) being immediate. So we may assume that \( n \geq 1 \). Let \( R_0 = F[[t_1, \ldots, t_{n-1}]] \) and \( R = R_0[[t_n]] = F[[t_1, \ldots, t_n]] \). By induction \( R_0 \) is a UFD, hence so is the polynomial ring \( R_0[t_n] \) (by Theorem 34.8). As \( R \) is Noetherian by Example 29.13(2), every nonzero nonunit in \( R \) is a product of irreducibles by Theorem 29.14. Let \( f \in R \) be irreducible. Therefore, it suffices to show that \( f \) is a prime element of \( R \). Suppose that \( f \mid rs \) in \( R \) with \( r \) and \( s \) in \( R \). Write \( fg = rs \) with \( g \in R \). By Lemma 36.3, there exists a ring automorphism \( \sigma \) of \( R \) fixing \( R_0[[t_n]] \) and satisfying \( \sigma(f), \sigma(g), \sigma(r), \) and \( \sigma(s) \) are all regular in \( t_n \). Therefore, we may assume that \( f, g, r, s \) are all regular in \( t_n \). By Corollary 36.6, we have \( f^*g^* = r^*s^* \), so we may further assume that \( f, g, r, s \) are all pseudo-polynomials in \( R_0[t_n] \). Since \( R_0[t_n] \) is a UFD, we are reduced to showing:
Claim. \( f \) is irreducible in \( R_0[t_n] \).

If \( f \in R_0[t_n] \) is reducible, then \( f = f_1 f_2 \) with \( f_1, f_2 \in R_0[t_n] \) satisfying \( \text{deg}_{t_n} f_i < \text{deg}_{t_n} f \) for \( i = 1, 2 \). (Why?) But \( f \) is irreducible in \( R \), so either \( f_1 \) or \( f_2 \) lies in \( R^\times \), say \( f_1 \). Then

\[
0 = f_1^{-1} f - f_2 \quad 0 = 0 \cdot f - 0
\]

in \( R \). By the uniqueness statement in the Weierstraß Preparation Theorem, \( f_1^{-1} = 0 \), which is impossible. This proves the claim and finishes the proof of the theorem. \( \square \)

Remark 36.10. In general, if \( R \) is a UFD, it does not follow that \( R[[t]] \) is a UFD. For example, it can be shown that \( R = \mathbb{F}[t_1, t_2, t_3]/(t_1^2 + t_2^3 + t_3^7) \) is a UFD but \( R[[t]] \) is not.

Exercises 36.11.
1. Prove that the map \( \sigma \) in the proof of Lemma 36.3 is well-defined.
2. Let \( R = \mathbb{F}[[t_1, \ldots, t_n]] \) with \( \mathbb{F} \) a field. Show that \( \mathfrak{m} = (t_1, \ldots, t_n) \) is the unique maximal ideal in \( R \) and \( \mathfrak{m}/\mathfrak{m}^2 \) is a vector space over \( \mathbb{F} \) of dimension \( n \).
Part 4

Module Theory
CHAPTER VIII

Modules

In this chapter, we study the basic result about modules. A module is a generalization of abelian groups and vector spaces. It is an object satisfying all the axioms of a vector space, except that the scalars are allowed to come from any fixed ring instead of from a field. The basic problem in module theory is given a ring with nice properties, classify modules over this ring up to isomorphism. Of course, this meets with relatively little success for all modules, so one looks at classes of modules over the ring.

In the first section, after giving examples of modules over a given ring, we establish the analogue of the isomorphism theorems and correspondence principle. Modules are nicer, in that we do not need a stronger property than submodules to state and prove these results compared to those needing normal subgroups and ideals in group and ring theory respectively. In particular, the quotient module of any given module and a submodule always exists.

We also study the special class of free modules. This class of modules generalizes the notion of vector space. More precisely free module are those that have bases, i.e., linearly independent spanning sets. This leads to the fact that dimension (called the rank) of a free module over a commutative ring makes sense, i.e., we can attach a unique integer to a finitely generated free module over a commutative ring. [Cf. the dimension of a finite dimensional vector space.]

37. Basic Properties of Modules

In this section, we introduce the concept of a module over a ring. As mentioned above, this generalizes the definition of vector spaces, by allowing scalars to come from an arbitrary ring rather than just a field. This difference, however, makes much that one expects in vector space theory to be false in this general case. One of the main reasons for this is that the concept of linear independence becomes a rare phenomenon. We begin with the formal definition.
Definition 37.1. Let $R$ be a ring. A (left) $R$-module is an additive group $M$ together with a map, $\cdot : R \times M \to M$ called scalar multiplication, which we write as $r \cdot x$ for $\cdot (r, x)$, satisfying for all $r, s \in R$ and $x, y \in M$:

(i) $r \cdot (x + y) = r \cdot x + r \cdot y$.
(ii) $(r + s) \cdot x = r \cdot x + s \cdot x$.
(iii) $r \cdot (s \cdot x) = (r \cdot s) \cdot x$.
(iv) $1 \cdot x = x$.

We call the map $\cdot : R \times M \to M$ an $R$-action. Note that there are two $\cdot$'s in the properties above, the $R$-action of $R$ on $M$ and the ring multiplication on $R$. Note that (i) and (ii) are the analogues of the distributive law and (iii) the analogue for the associative law for rings. Condition (iv) says the identity of the ring acts like an identity. You should also compare this to $G$-sets where $G$ is an additive group (or a monoid). We drop the $\cdot$ of the $R$-action whenever it is clear.

A subgroup $N$ in $M$ is called a submodule if it is an $R$-module under $\cdot|_{R \times N}$.

Remarks 37.2. 1. One defines right $R$-modules in the obvious way.

2. As usual, a submodule of an $R$-module $M$ should be a subset of $M$ that is an $R$-module under the inclusion map. It will be so after we define the appropriate notion of homomorphism. Of course, you should already know what it is, as we are generalizing vector space theory.

3. Some authors call $M$ an $R$-module without assuming (iv) and call an $M$ also satisfying (iv) a unitary $R$-module. (If $R$ is a rng, we will not assume (iv) in our definition.)

4. Let $M$ be an $R$-module, $x$ an element in $M$, and $r$ an element in $R$. Then

(a) $0_R x = 0_M$. (We write $0x = 0$.)
(b) $r0_M = 0_M$. (We write $r0 = 0$.)
(c) $(-r)x = r(-x) = -(rx)$.
(d) $m(rx) = r(mx)$ for all integers $m$.

We leave this an an exercise.

As usual we start with a number of examples.

Examples 37.3. Let $R$ be a ring and $M$ an $R$-module.

1. If $R$ is a field (or a division ring), then an $R$-module is the same as a vector space over $R$. 

2. If $R$ is the ring of integers, then a $\mathbb{Z}$-module is just an additive group. In particular, module theory generalizes both vector space theory and abelian group theory.

3. $0 := \{0\}$ and $M$ are submodules of $M$.

4. $R$ is an $R$-module under the addition in $R$ with the $R$-action given by the multiplication in $R$.

5. Let $\mathfrak{A}$ be a left ideal in $R$. Then $\mathfrak{A} + \mathfrak{A} \subset \mathfrak{A}$ and $R\mathfrak{A} \subset \mathfrak{A}$, so $\mathfrak{A}$ is a submodule of $R$. Indeed, a submodule of $R$ is the same thing as a left ideal. So the concept of an $R$-module also generalizes the notion of left ideal.

6. $R[t]$ is an $R$-module. More generally, if $R$ is a subring of $S$ then $S$ is an $R$-module by restricting the ring multiplication on $S$ to scalar multiplication by $R$, i.e., restricting $\cdot_{|S \times S}$ to $\cdot_{|R \times S}$.

7. Let $S$ be a nonempty set and $X$ an $R$-module. Set

$$M := \{f : S \to X \mid f \text{ a map}\},$$

then $M$ is an $R$-module under the following operations: If $f$ and $g$ lie in $M$ and $r$ in $R$, then for all $s$ in $S$ define

$$(f + g)(s) := f(s) + g(s)$$

and $0_M$ defined by $0_M(s) = 0_X$.

8. Let $N$ be a submodule of $M$. Then the factor group $M/N$ ($N \triangleleft M$ is automatic as $M$ is abelian) becomes an $R$-module by

$$\cdot : R \times M/N \to M/N \text{ defined by } r(m + N) = rm + N$$

for all $r \in R$ and for all $m \in M$. It is called the quotient or factor module of $M$ by $N$.

9. Let $N$ be an abelian group. Write it additively, i.e., $N$ is a $\mathbb{Z}$-module. Then

$$\text{End}(N) = \text{End}_\mathbb{Z}(N) := \{ \varphi : N \to N \mid \varphi \text{ a group homomorphism} \}$$

is a ring via

$$(i) \quad (\varphi + \psi)(x) = \varphi(x) + \psi(x) \text{ (the usual addition of functions into an additive group)}$$

$$(ii) \quad \varphi \circ \psi(x) = \varphi(\psi(x)) \text{ (composition of functions)}.$$ 

for all $\varphi$ and $\psi$ in $\text{End}(N)$ and for all $x$ in $N$ (as $nx = \underbrace{x + \cdots + x}_n$ for all positive integers $n$); and $N$ is an $\text{End}(N)$-module by evaluation, i.e., $\varphi \cdot x := \varphi(x)$ for all $\varphi \in \text{End}(N)$ and $x \in N$. 

10. Let $F$ be a field, $V$ a vector space over $F$. Then

$$\text{End}_F(V) = \{ T : V \to V \mid T \text{ a linear operator} \}$$

is a subring of $\text{End}_Z(V)$ and $V$ is an $\text{End}_F(V)$-module via evaluation, i.e., $T \cdot x = T(x)$ for all $T \in \text{End}_F(V)$ and $x \in V$.

[Can you define the appropriate notion of homomorphism $f : M \to N$ between modules, then replace $\text{End}_F(V)$ by the appropriate $\text{End}_R(M)$?]

11. Let $\varphi : R \to S$ be a ring homomorphism and $N$ an $S$-module. Then $N$ becomes an $R$-module via the pullback defined by

$$\cdot : R \times N \to N \text{ given by } r \cdot x = \varphi(r) \cdot x$$

for all $r$ in $R$ and $x$ in $N$.

12. Let $V$ be a vector space over the field $F$ and $T \in \text{End}_F(V)$ a fixed linear operator (= endomorphism). Define a ring homomorphism

$$\varphi : F[t] \to \text{End}_F(V) \text{ given by the evaluation } f \to f(T).$$

So we have $\sum a_i t^i \mapsto \sum a_i T^i$. [What is ker $\varphi$? What can you say if $V$ is a finite dimensional vector space over $F$?] Then $V$ is an $F[t]$-module via the pullback $f \cdot v := f(T)(v)$ for all $f \in F[t]$ and for all $v \in V$. Let $W$ be a subspace of $V$. Then $W$ is a submodule of the $F[t]$-module $V$ if and only if $W$ is a $T$-invariant subspace, i.e., $T(W) \subseteq W$. This is an example that will be used in Section 44.

13. Let $\mathfrak{a} \subset R$ be a left ideal. Then

$$\mathfrak{a} M := \{ \sum a_x x \mid a_x \in \mathfrak{a} \text{ for all } x \in M \} \subset M$$

is an $R$-module. Note that an infinite sum of elements in an additive group does not make sense, so when we write $\sum_{x \in M} a_x x$, we must have $a_x$ is zero for almost all $x$, i.e., $a_x \neq 0$ for only finitely many $x$.

14. Let $M$ be an $R$-module and $M_i$ submodules of $M$ for $i \in I$. Then

(i) $\bigcap_{i} M_i$ and

(ii) $\sum_{i} M_i := \{ \sum m_i \mid m_i \in M_i, m_i = 0 \text{ for almost all } i \in I \}$

are submodules of $M$. If the $M_i$ also satisfy $M_i \cap \sum_{j \notin i} M_j = 0$ for all $i \in I$, we call $\sum_{i} M_i$ the internal direct sum of the $M_i$ and write it as $\bigoplus_{i} M_i$.

15. Let $M_i$, $i \in I$, be $R$-modules. Recall $(m_i)_I$ is the notation for elements in the cartesian product $\otimes_I M_i$. Set
37. BASIC PROPERTIES OF MODULES

(i) \( \prod_I M_i := \{ (m_i)_I \mid m_i \in M_i, m_i = 0 \text{ for almost all } i \in I \} \), an \( R \)-module under componentwise operations. It is called the \( (\text{external}) \) direct sum of the \( M_i \).

(ii) \( \coprod_I M_i := \{ (m_i)_I \mid m_i \in M_i \} \), an \( R \)-module under componentwise operations. It is called the \( (\text{external}) \) direct product of the \( M_i \).

We have \( \coprod_I M_i \) is a submodule of \( \prod_I M_i \) with \( \coprod_I M_i = \prod_I M_i \) if and only if \( I \) is finite or almost all \( M_i = 0 \).

16. Let \( X \) be a subset of \( M \). Define

\[ \langle X \rangle := \sum X Rx, \]

called the submodule \( \text{generated} \) by \( X \). Note that, as for vector spaces, this is equivalent to \( \langle X \rangle = \bigcap_{X \subseteq N \subseteq M} N \), where \( N \) runs over all submodules of \( M \) containing \( X \). If there exists a finite subset \( Y \) in \( M \) such that \( M = \langle Y \rangle \) we say that \( M \) is \( \text{finitely generated} \) or simply \( \text{fg} \). If \( Y = \{ y_1, \ldots, y_n \} \), we write \( \langle y_1, \ldots, y_n \rangle \) for \( \langle Y \rangle \). If there exists a \( y \in M \) such that \( M = \langle y \rangle \), we say that \( M \) is \( R \)-cyclic. (So a cyclic group is the same thing as a cyclic \( \mathbb{Z} \)-module.)

**Definition 37.4.** A map \( f : M \to N \) of \( R \)-modules is called an \( R \)-homomorphism (respectively, \( R \)-monomorphism, \( R \)-epimorphism, \( R \)-isomorphism) if \( f \) is \( R \)-linear, i.e.,

\[ f(rx + y) = rf(x) + f(y) \]

for all \( r \in R \) and \( x, y \in M \) (respectively, and injective, and surjective, and bijective with inverse \( R \)-linear). As usual we also use the abbreviations of monic and epic. Also, as one would surmise, \( f \) is an \( R \)-isomorphism if and only if it is a bijective \( R \)-homomorphism.

We say that two \( R \)-modules \( M \) and \( N \) are \( R \)-isomorphic and write \( M \cong N \) if there exists an \( R \)-isomorphism \( g : M \to N \), which we also write as \( g : M \xrightarrow{\sim} N \).

**Remark 37.5.** As is true with groups and rings an \( R \)-homomorphism \( f : M \to N \) is an \( R \)-monomorphism if and only if given any \( R \)-homomorphisms \( g_1, g_2 : L \to M \) with compositions satisfying \( f \circ g_1 = f \circ g_2 \), then \( g_1 = g_2 \); and as with groups (but not rings) \( f \) is an \( R \)-epimorphism if and only if given any \( R \)-homomorphisms \( h_1, h_2 : N \to L \) with compositions satisfying \( h_1 \circ f = h_2 \circ f \), then \( h_1 = h_2 \). (Cf. Exercise 1.12(6) and (7).)

**Examples 37.6.** 1. A \( \mathbb{Z} \)-homomorphism is the same thing as an abelian group homomorphism.
2. A linear transformation of $F$-vector spaces is the same thing as an $F$-homomorphism of $F$-modules.

3. Let $M$ and $N$ be $R$-modules with $N \subset M$. Then $N$ is a submodule of $M$ if and only if the inclusion is an $R$-homomorphism.

4. Let $R$ be a commutative ring and $r$ an element of $R$. Then
   \[ \lambda_r : R \to R \text{ given by } x \mapsto rx \]
is an $R$-homomorphism (but not a ring homomorphism). More generally, if $M$ is an $R$-module, then
   \[ \lambda_r : M \to M \text{ given by } x \mapsto rx \]
is an $R$-homomorphism.

5. Let $N$ be a submodule of $M$. Then the canonical map
   \[ \pi : M \to M/N \text{ given by } x \mapsto x + N \]
is an $R$-epimorphism.

6. Let $R$ be a ring and $\mathfrak{a} < R$ be a (2-sided) ideal, \( \pi : R \to R/\mathfrak{a} \)
   the canonical ring epimorphism. If $M$ is an $R$-module then the $R$-module $M/\mathfrak{a}M$ is also an $R/\mathfrak{a}$-module by \( \pi(x + \mathfrak{a}M) = rx + \mathfrak{a}M \)
   for all $r \in R$ and $x \in M$, so the $R$- and $(R/\mathfrak{a})$-actions on $M/\mathfrak{a}M$ are compatible, i.e., $r(x + \mathfrak{a}M) = \pi(x + \mathfrak{a}M)$.

7. Let $f : M \to N$ be an $R$-homomorphism of $R$-modules. Then the additive subgroups $\ker f \subset M$ and $\im f \subset N$ are submodules. More generally, if $N'$ is a submodule of $N$, then $\ker f \subset f^{-1}(N') \subset M$ is a submodule and if $M'$ is a submodule of $M$, then $f(M')$ is a submodule of $N$. This makes the study of $R$-modules and $R$-homomorphisms very nice. We do not need to worry about some of the special properties needed with groups and rings.

8. If $M$ is an $R$-module, then
   \[ \text{End}_R(M) := \{ f : M \to M \mid f \text{ an } R\text{-homomorphism} \} \]
is a ring under the + of functions and composition (as asked for before), called the endomorphism ring of the $R$-module $M$. Elements of $\text{End}_R(M)$ are called $R$-endomorphisms. Its unit group is $\text{Aut}_R(M)$, the group of $R$-automorphisms of $M$, where an $R$-automorphism is an $R$-isomorphism $M \to M$. The ring $\text{End}(M) = \text{End}_Z(M)$ contains $\text{End}_R(M)$ as a subring. Note that if $R$ is a field, then a linear operator $T : V \to V$ of a vector space $V$ over $R$ is the same thing as an $R$-endomorphism of $V$. 
9. If $R$ is a commutative ring and $M$, $N$ are $R$-modules, then

$$\text{Hom}_R(M, N) := \{ f : M \to N \mid f \text{ an } R\text{-homomorphism} \}$$

is an $R$-module with the usual $+$ for functions and the $R$-action $\cdot$ given by $r \cdot f : x \mapsto rf(x)$. [If $R$ is not commutative, $\text{Hom}_R(M, N)$ is not an $R$-module and we usually just view it as an abelian group. Is it an $S$-module for some subrings $S$ of $R$?]

We have analogous isomorphism theorems and a correspondence theorem for modules as we had before with the same proof. We state them for convenience.

**Theorem 37.7.** (First Isomorphism Theorem) Let $f : M \to N$ be an $R$-homomorphism. Then we have a commutative diagram of $R$-modules and $R$-homomorphisms

$$
\begin{array}{ccc}
M & \xrightarrow{f} & N \\
\downarrow & & \uparrow \text{inc} \\
M/\ker f & \xrightarrow{\overline{f}} & \text{im } f
\end{array}
$$

with $\overline{f}$, the canonical map, an $R$-epimorphism, $\overline{f}$, the map given by $x \mapsto f(x)$, an $R$-isomorphism between $M/\ker f$ and $\text{im } f$, and the inclusion map $\text{inc}$ an $R$-monomorphism.

**Theorem 37.8.** (Second Isomorphism Theorem) Let $M$ be an $R$-module and $A$ and $N$ submodules of $M$. Then $A/(A \cap N) \cong (A + N)/N$.

**Theorem 37.9.** (Third Isomorphism Theorem) Let $M$ be an $R$-module with submodules $A$ and $N$ satisfying $A \subset N$. Then we have $M/N \cong (M/A)/(N/A)$.

**Theorem 37.10.** (Correspondence Principle) Let $f : M \to N$ be an $R$-epimorphism of $R$-modules. Then

$$\{ A \mid \ker f \subset A \subset M \text{ a submodule} \} \longrightarrow \{ B \mid B \subset N \text{ a submodule} \}$$
given by $A \mapsto f(A)$ is an order preserving bijection.

Next we introduce the generalization of a zero divisor in a ring.

**Definition 37.11.** Let $M$ be an $R$-module and $m$ an element of $M$. The **annihilator** of $m$ is defined to be the set

$$\text{ann}_R m := \{ r \in R \mid rm = 0 \}.$$ 

More generally, we shall let

$$\text{ann}_R(M) := \{ r \in R \mid rm = 0 \text{ for all } m \in M \}.$$
denote the annihilator of $M$.

The annihilator has the following properties:

**Lemma 37.12.** Let $M$ be an $R$-module and $m$, $m'$ elements in $M$. Then

1. $\text{ann}_R(M) \subset R$ is a left ideal.
2. $\rho_m : R \to M$ given by $r \mapsto rm$ is an $R$-homomorphism and satisfies $\ker \rho_m = \text{ann}_R m$.
3. If $\rho_m$ is the $R$-homomorphism in (2), then $\rho_m$ induces an $R$-isomorphism $\overline{\rho}_m : R/\text{ann}_R m \sim \to Rm$.
4. If $R$ is a commutative ring and $Rm \subset Rm'$, then $\text{ann}_R m \supset \text{ann}_R m'$. In particular, $\text{ann}_R m$ is independent of the generator for $Rm$.

**Proof.** (1) and (2) follow immediately. (3) follows from the First Isomorphism Theorem for modules. (4): Suppose that we know that $m = am'$ for some $a \in R$. If $rm' = 0$, then $ram' = arm' = 0$. □

We can now characterize cyclic $R$-modules.

**Corollary 37.13.** Let $M$ be an $R$-module. Then $M$ is a cyclic $R$-module if and only if there exists a left ideal $\mathfrak{a}$ in $R$ satisfying $M \cong R/\mathfrak{a}$.

**Proof.** $R/\mathfrak{a} = (1 + \mathfrak{a}) = R(1 + \mathfrak{a})$ is cyclic, so this follows by the lemma. □

If $R$ is a commutative ring, we can say more.

**Proposition 37.14.** Let $R$ be a commutative ring, $M$ a cyclic $R$-module. Then there exists a unique ideal $\mathfrak{a}$ in $R$ such that $M \cong R/\mathfrak{a}$. Moreover, if $M = Rm$, then $\mathfrak{a} = \text{ann}_R m$. In particular, if $R$ is a PID, there exists an element $a$ in $R$, unique up to units, satisfying $M \cong R/(a)$, and, if in addition, $M = Rm$, then $(a) = \text{ann}_R m$.

**Proof.** We already know if $M = Rm$, then $M \cong R/\text{ann}_R m$. Suppose that we have an $R$-isomorphism $f : R/\mathfrak{a} \to M$ for some ideal $\mathfrak{a}$ of $R$. Let $m = f(1 + \mathfrak{a})$. We have

$$f(r + \mathfrak{a}) = f(r(1 + \mathfrak{a})) = rf(1 + \mathfrak{a}) = rm$$
for all \( r \in R \). In particular, \( M = Rm \). Let \( \rho_m : R \rightarrow M \) be the \( R \)-epimorphism given by \( r \mapsto rm \). Then we have a commutative diagram

\[
\begin{array}{c}
R \\
\downarrow \quad \downarrow f \\
R/\mathfrak{A}.
\end{array}
\]

Thus \( \text{ann}_R m = \ker \rho_m = \ker(f \circ -) = \ker - = \mathfrak{A} \),

as \( f \) is monic and \( f(\tau) = \rho_m(r) \), i.e., \( f = \overline{\rho_m} \). By Lemma 37.12(4), \( \text{ann}_R m \) is independent of the generator of \( M \) when \( R \) is commutative. If follows that \( \mathfrak{A} \) is unique. Lemma 37.12(4) also implies the statements about PIDs. \( \Box \)

We now set up a lot of useful notation. We write maps \( f : A \rightarrow B \) as \( A \xrightarrow{f} B \) as we have done in commutative diagrams.

**Definition 37.15.** Let \( A \xrightarrow{f} B \) and \( B \xrightarrow{g} C \) be \( R \)-homomorphisms of \( R \)-modules. We say that

\( A \xrightarrow{f} B \xrightarrow{g} C \)

is a zero sequence if \( g \circ f = 0 \), equivalently, \( \text{im } f \subseteq \ker g \) and is an exact sequence if \( \text{im } f = \ker g \). A longer sequence

\[
\cdots \xrightarrow{f_{n+2}} A_{n+1} \xrightarrow{f_{n+1}} A_n \xrightarrow{f_n} A_{n-1} \xrightarrow{f_{n-1}} \cdots
\]

of \( R \)-modules and \( R \)-homomorphisms is called a zero sequence or a chain complex (respectively, an acyclic complex) if \( f_n f_{n+1} = 0 \) (respectively, \( \text{im } f_{n+1} = \ker f_n \)) for all \( n, n+1 \) occurring. The special case of an exact sequence of \( R \)-modules and \( R \)-homomorphisms of the form

\[
0 \rightarrow A \xrightarrow{f} B \xrightarrow{g} C \rightarrow 0
\]

is called a short exact sequence.

**Examples 37.16.** Let \( A \xrightarrow{f} B \) and \( B \xrightarrow{g} C \) be \( R \)-homomorphisms of \( R \)-modules.

1. The \( R \)-homomorphism \( f \) is monic if and only if \( 0 \rightarrow A \xrightarrow{f} B \) is exact. [Note that \( 0 \rightarrow A \) must take \( 0 \mapsto 0_A \) as it is an \( R \)-homomorphism.]
2. The \( R \)-homomorphism \( g \) is epic if and only if \( B \xrightarrow{g} C \rightarrow 0 \) is exact.
3. If \( 0 \rightarrow A \xrightarrow{f} B \xrightarrow{g} C \rightarrow 0 \) is a short exact exact sequence, then \( f \) is monic, \( \text{im } f = \ker g \), and \( g \) is epic.
4. The sequence $0 \to \ker f \to A \xrightarrow{f} B$ is exact (where the second map is the inclusion) and $0 \to \ker f \to A \xrightarrow{f} \text{im } f \to 0$ is a short exact sequence (writing the same $f$ after changing the target).

5. The $R$-homomorphism $f$ induces an $R$-monomorphism $A/\ker f \xrightarrow{\overline{f}} B$ by the First Isomorphism Theorem. Therefore, the First Isomorphism Theorem is just that we have a commutative diagram

\[
\begin{array}{ccccccccc}
0 & \to & \ker f & \to & A & \xrightarrow{\overline{f}} & A/\ker f & \to & 0 \\
& & \parallel & & \parallel & & \downarrow \overline{f} \\
0 & \to & \ker f & \to & A & \xrightarrow{f} & \text{im } f & \to & 0.
\end{array}
\]

with exact rows and with the right hand vertical arrow $\overline{f}$ an $R$-isomorphism. Can you write the Third Isomorphism Theorem using such sequences?

6. As the factor $R$-module of $M$ by $N$ exists for any submodule $N$ of an $R$-module $M$, besides the image and kernel of the $R$-homomorphism $B \xrightarrow{g} C$, we can define another $R$-module, the cokernel of $g$ defined by $\text{coker } g := C/\text{im } g$. Then we have a short exact sequence

\[
0 \to \text{im } g \to C \xrightarrow{\pi} \text{coker } g \to 0,
\]

where the second map is the inclusion. So we have

(a) $g$ is injective if and only if $\ker g = 0$.

(b) $g$ is surjective if and only if $\text{coker } g = 0$.

7. Let $A$ and $B$ be $R$-modules. Then the sequence

\[
0 \to A \xrightarrow{\iota_A} A \coprod B \xrightarrow{\pi_B} B \to 0
\]

with $\iota_A(a) = (a, 0)$ and $\pi_B(a, b) = b$ is a short exact sequence.

8. Let $A_i$, $B_i$, and $C_i$ be $R$-modules for $i \in I$. If

\[
0 \to A_i \xrightarrow{f_i} B_i \xrightarrow{g_i} C_i \to 0
\]

is a short exact sequence for all $i \in I$, then

\[
0 \to \coprod_I A_i \xrightarrow{\coprod_I f_i} \coprod_I B_i \xrightarrow{\coprod_I g_i} \coprod_I C_i \to 0
\]

and

\[
0 \to \prod_I A_i \xrightarrow{\prod_I f_i} \prod_I B_i \xrightarrow{\prod_I g_i} \prod_I C_i \to 0
\]

are exact. What are the maps $\coprod_I f_i$, $\coprod_I g_i$, and $\prod_I f_i$, $\prod_I g_i$?
Lemma 37.17. (Five Lemma) Suppose the following is a commutative diagram of $R$-modules and $R$-homomorphisms with exact rows:

\[
\begin{array}{ccccccccc}
0 & \longrightarrow & A & \xrightarrow{f} & B & \xrightarrow{g} & C & \longrightarrow & 0 \\
& h_A & & h_B & & h_C & & & \\
0 & \longrightarrow & A' & \xrightarrow{f'} & B' & \xrightarrow{g'} & C' & \longrightarrow & 0.
\end{array}
\]

If two of $h_A$, $h_B$, $h_C$ are $R$-isomorphisms, then they all are.

The proof of this lemma is called diagram chasing, and we leave it as an exercise.

Remark 37.18. A chain complex is usually written

\[
\ldots \rightarrow d_{n+2}A_{n+1} \xrightarrow{d_{n+1}} A_n \xrightarrow{d_n} A_{n-1} \xrightarrow{d_{n-2}} \ldots
\]

and the maps $d_i$ are called differentials. One usually denotes such a chain complex by $(A_*, d_*)$. If $(A_*, d_*)$ is a chain complex, then we define the quotient of $R$-modules

\[
H_n(A) := \ker d_n / \text{im} d_{n+1}
\]

called the $n$th homology of $(A_*, d_*)$. If $(A_*, d_*)$ is a chain complex, the set $A_n$ is called the set of $n$-chains. We call

\[
Z_n(A) := \ker d_n \quad \text{the set of } n\text{-cycles of } A_n
\]

\[
B_n(A) := \text{im} d_{n+1} \quad \text{the set of } n\text{-boundaries of } A_n.
\]

So $H_n(A) = Z_n(A)/B_n(A)$ and every cohomology class in $H_n(A)$ is represented by an $n$-cycle. This agrees with its historical topological interpretation.

If the indices on the $d_n$ go up, i.e.,

\[
\ldots \rightarrow d_{n-2}A_{n-1} \xrightarrow{d_{n-1}} A_n \xrightarrow{d_n} A_{n+1} \xrightarrow{d_{n+1}} \ldots,
\]

then $(A_*, d_*)$ is called a cochain complex with differentials written as $d^n$. We use the notation $(A_*, d^n)$ for a cochain complex. Analogously, the elements in $A_n$ are called $n$-cochains and

\[
Z^n(A) := \ker d^n \quad \text{is called the set of } n\text{-cocycles of } A_n
\]

\[
B^n(A) := \text{im} d^{n-1} \quad \text{is called the set of } n\text{-coboundaries of } A_n
\]

\[
H^n(A) := \ker d^n / \text{im} d^{n-1} = Z^n(A)/B^n(A).
\]

Therefore, a (co)chain complex is exact, i.e., acyclic, if and only if its (co)homology is trivial.
A chain homomorphism (or chain homomorphism of degree zero) of chain complexes \( f : (A_*, d_*) \to (A'_*, d'_*) \) consists of \( R \)-modules and \( R \)-homomorphisms is \( f_n : A_n \to A'_n \) for all \( n \) satisfying

\[
\begin{array}{ccc}
A_{n+1} & \xrightarrow{f_{n+1}} & A'_{n+1} \\
\downarrow{d_{n+1}} & & \downarrow{d'_{n+1}} \\
A_n & \xrightarrow{f_n} & A'_n
\end{array}
\]

commutes for all \( n \). The obvious definition is given for cochain homomorphisms.

The study of (co)homology, in its general context, is a major subject in mathematics. We shall leave proofs and statements of some of the basic results in the exercises. These exercises may be safely omitted until the reader is more comfortable with module theory.

**Exercises 37.19.**

1. Prove Remark 37.2(4)
2. Prove Remark 37.5.
3. Show an \( R \)-homomorphism \( f : M \to N \) is an isomorphism if and only if \( f \) is bijective.
4. Let \( M \) be a simple \( R \)-module, i.e., \( M \) has no proper nonzero submodules. Prove that \( \text{End}_R(M) \) is a division ring.
5. Let \( M \) be an \( R \)-module and \( M_i \) submodules of \( M \) for \( i \in I \). Show that there is always an \( R \)-homomorphism \( f : \bigsqcup_i M_i \to \sum_i M_i \) and this homomorphism is an \( R \)-isomorphism if and only if \( \sum_i M_i \) is an internal direct sum.
6. Let \( M \) be an \( R \)-module and \( M_1, \ldots, M_n \) be \( R \)-submodules of \( M \). Show that \( M = \bigoplus_{i=1}^n M_i \) if and only if there exist \( R \)-homomorphisms \( \iota_i : M_i \to M \) and \( \pi_i : M \to M_i \), \( i = 1, \ldots, n \), satisfying all of the following:
   
   \[ (i) \quad \pi_i \iota_i = 1_{M_i} \text{ for } i = 1, \ldots, n. \]
   
   \[ (ii) \quad \pi_j \iota_i = 0 \text{ for } i \neq j. \]
   
   \[ (iii) \quad \iota_1 \pi_1 + \cdots + \iota_n \pi_n = 1_M. \]
7. Let \( m \) be a positive integer. Determine the abelian groups (\( \mathbb{Z} \)-modules) \( \text{Hom}_\mathbb{Z}(\mathbb{Z}, \mathbb{Z}/m\mathbb{Z}) \) and \( \text{Hom}_\mathbb{Z}(\mathbb{Z}/m\mathbb{Z}, \mathbb{Z}) \) up to isomorphism.
8. Let \( m \) and \( n \) be positive integers with greatest common divisor \( d \). Show \( \text{Hom}_\mathbb{Z}(\mathbb{Z}/m\mathbb{Z}, \mathbb{Z}/n\mathbb{Z}) \cong \mathbb{Z}/d\mathbb{Z} \).
9. (Universal Property of Direct Sums) Let \( M_i, i \in I \), be \( R \)-modules. Show that the (external) direct sum \( \bigsqcup_i M_i \) satisfies the following universal property relative to the \( R \)-homomorphisms \( \iota_j : M_j \to \bigsqcup_i M_i \).
defined by \( m \mapsto \{ \delta_{ij} m \}_{i \in I} \) for all \( j \in I \): Given an \( R \)-module \( M \) and \( R \)-homomorphisms \( f_j : M_j \to M \) for all \( j \in I \), there exist a unique \( R \)-homomorphism \( g : \prod_i M_i \to M \) satisfying \( f_j = g \circ \iota_j \) for all \( j \in I \).

10. (Universal Property of Direct Products) Let \( M_i, i \in I \), be \( R \)-modules. Show that the (external) direct product \( \prod_i M_i \) satisfies the following universal property relative to the \( R \)-homomorphisms \( \pi_j : \prod_i M_i \to M_j \) defined by \( \{ m_i \} \mapsto m_j \) for all \( j \in I \): Given an \( R \)-module \( M \) and \( R \)-homomorphisms \( g_j : M \to M_j \) for all \( j \in I \), there exist a unique \( R \)-homomorphism \( h : M \to \prod_i M \) satisfying \( g_j = \pi_j \circ h \) for all \( j \in I \).

11. Write the Third Isomorphism Theorem using exact sequences as in Example 37.16(5).

12. A short exact sequence of \( R \)-modules

\[
0 \to A \xrightarrow{f} B \xrightarrow{g} C \to 0
\]

is called \textit{split} if one of the following three equivalent conditions holds:

(i) There exists an \( R \)-homomorphism \( f' : B \to A \) such that \( f'f = 1_A \). We say that \( f \) is a \textit{split monomorphism}.

(ii) \( \text{im } f \) is a \textit{direct summand} of \( B \), i.e., \( B = \text{im } f \oplus D \) some \( R \)-module \( D \).

(iii) There exists an \( R \)-homomorphism \( g' : C \to B \) such that \( gg' = 1_C \). We say that \( g \) is a \textit{split epimorphism}.

Prove that these conditions are equivalent.

13. Prove the Five Lemma 37.17.

14. Prove the full version of the Five Lemma: Suppose that

\[
A \xrightarrow{\alpha} B \xrightarrow{\beta} C \xrightarrow{\gamma} D \xrightarrow{\delta} E
\]

is a commutative diagram of \( R \)-modules and \( R \)-homomorphisms with exact rows. Then show the following:

(i) If \( h_A \) is an epimorphism and \( h_B \) and \( h_D \) are monomorphisms, then \( h_C \) is a monomorphism.

(ii) If \( h_E \) is a monomorphism and \( h_B \) and \( h_D \) are epimorphisms, then \( h_C \) is an epimorphism.

is a commutative diagram of \( R \)-modules and \( R \)-homomorphisms with exact rows.

**Homology Theory.** The following exercises develop the basics of homology theory.
15. (Snake Lemma) Let

\[
\begin{array}{ccc}
A & \xrightarrow{f} & B \\
\downarrow{\alpha} & & \downarrow{\beta} \\
0 & \xrightarrow{f'} & A'
\end{array}
\quad
\begin{array}{ccc}
& & C \\
\downarrow{\gamma} & & \\
& & \\
\downarrow{\gamma'} & & \downarrow{\gamma'} \\
0 & \xrightarrow{g'} & C'
\end{array}
\quad
\begin{array}{ccc}
& & 0 \\
\end{array}
\]

be a commutative diagram of \(R\)-modules and \(R\)-homomorphisms with exact rows. Prove that

\[\partial : \ker \gamma \to \coker \alpha \text{ defined by } a \mapsto f' - 1\beta g^{-1}a + \im \alpha\]
is a well-defined \(R\)-homomorphism and the following sequence is exact:

\[
\ker \alpha \xrightarrow{f_{|\ker \alpha}} \ker \beta \xrightarrow{g_{|\ker \beta}} \ker \gamma \xrightarrow{\partial} \coker \alpha \xrightarrow{\overline{f}} \coker \beta \xrightarrow{\overline{g}} \coker \gamma
\]

Moreover, in this diagram show that \(f_{|\ker \alpha}\) is a monomorphism if \(f\) is a monomorphism and \(\overline{g}\) is an epimorphism if \(g\) is an epimorphism.

16. Let

\[
\begin{array}{ccccccccc}
0 & \rightarrow & A_{n+1} & \xrightarrow{f_{n+1}} & B_{n+1} & \xrightarrow{g_{n+1}} & C_{n+1} & \rightarrow & 0 \\
& & \downarrow{(d_A)_{n+1}} & & \downarrow{(d_B)_{n+1}} & & \downarrow{(d_C)_{n+1}} & & \\
0 & \rightarrow & A_n & \xrightarrow{f_n} & B_n & \xrightarrow{g_n} & C_n & \rightarrow & 0 \\
& & \downarrow{(d_A)_n} & & \downarrow{(d_B)_n} & & \downarrow{(d_C)_n} & & \\
0 & \rightarrow & A_{n-1} & \xrightarrow{f_{n-1}} & B_{n-1} & \xrightarrow{g_{n-1}} & C_{n-1} & \rightarrow & 0.
\end{array}
\]

be a commutative diagram of \(R\)-modules and \(R\)-homomorphisms having rows exact and verticals zero sequences. Show the following diagram is commutative with exact rows:

\[
\begin{array}{cccc}
\coker (d_A)_{n+1} & \xrightarrow{\overline{f}_{n+1}} & \coker (d_B)_{n+1} & \xrightarrow{\overline{g}_{n+1}} & \coker (d_C)_{n+1} & \rightarrow & 0 \\
\downarrow{\overline{d}_A} & & \downarrow{\overline{d}_B} & & \downarrow{\overline{d}_C} & & \\
0 & \rightarrow & \ker (d_A)_{n-1} & \xrightarrow{(f_{n-1})|_{\ker d_A}} & \ker (d_B)_{n-1} & \xrightarrow{(g_{n-1})|_{\ker d_B}} & \ker (d_C)_{n-1} & &
\end{array}
\]

where the maps \(\overline{f}, \overline{g}, \overline{d}_A, \overline{d}_B, \overline{d}_C\) are the induced maps.

17. (The Long Exact Sequence in Homology) Let

\[
0 \rightarrow (A_*, d_A) \xrightarrow{f_*} (B_*, d_B) \xrightarrow{g_*} (C_*, d_C) \rightarrow 0
\]
be a short exact sequence of chain complexes of $R$-homomorphisms and $R$-modules (where 0 is the trivial complex $(0,0)$) i.e.,

$$0 \to A_n \xrightarrow{f_n} B_n \xrightarrow{g_n} C_n \to 0$$

is a short exact sequence for all $n$. Show that there exists an $R$-homomorphism $\partial_{n+1} : H_{n+1}(C) \to H_n(A)$ for all $n$ (arising from the Snake Lemma) called the connecting homomorphism of the sequence such that we have a long exact sequence in homology

$$\cdots \to H_{n+1}(C) \xrightarrow{\partial_{n+1}} H_n(A) \xrightarrow{\bar{f}_n} H_n(B) \xrightarrow{\bar{g}_n} H_n(C) \xrightarrow{\partial_n} H_{n-1}(A) \to \cdots$$

where $\bar{f}_n$ and $\bar{g}_n$ are the induced maps for all $n$.

18. Let

$$0 \longrightarrow A_* \xrightarrow{f} B_* \xrightarrow{g} C_* \longrightarrow 0$$

be a commutative diagram of chain complexes of $R$-modules with exact rows. Show that there exists a commutative diagram with exact rows

$$\cdots \longrightarrow H_n(A) \xrightarrow{f_n} H_n(B) \xrightarrow{g_n} H_n(C) \xrightarrow{\partial_n} H_{n-1}(A) \longrightarrow \cdots$$

$$\quad \downarrow \pi_n \quad \downarrow \beta_n \quad \downarrow \gamma_n \quad \downarrow \alpha_{n-1}$$

$$\cdots \longrightarrow H_n(A') \xrightarrow{f'_n} H_n(B') \xrightarrow{g'_n} H_n(C') \xrightarrow{\partial'_n} H_{n-1}(A') \longrightarrow \cdots$$

where $\partial_n$ and $\partial'_n$ are the corresponding connecting homomorphisms and $\pi_n$, $\beta_n$, and $\gamma_n$ the induced maps for all $n$.

19. Let $f_*, g_* : (A_*, d_A) \to (B_*, d_B)$ be chain maps of chain complexes of $R$-modules. Let $s_n : A_n \to B_{n+1}$ be an $R$-homomorphism for all $n$ and $s_* = (s_n)$ (called a chain map of degree $+1$). We say that $f$ and $g$ are chain homotopic if there exists a chain map $s_* = (s_n)$ of degree $+1$ that satisfies $f_n - g_n = (d_B)_{n+1}s_n + s_{n-1}(d_A)_n$, i.e., we have a commutative diagram

$$\begin{array}{ccc}
A_{n+1} & \xrightarrow{(d_A)_{n+1}} & A_n & \xrightarrow{(d_A)_n} & A_{n-1} \\
\downarrow f_{n+1} & & \downarrow f_n & & \downarrow f_{n-1} \\
B_{n+1} & \xrightarrow{(d_B)_{n+1}} & B_n & \xrightarrow{(d_B)_n} & B_{n-1}
\end{array}$$
Show if \( f_* \) and \( g_* \) are chain homotopic, then the maps \( f_n = g_n : H_n(A) \to H_n(B) \) for all \( n \).

20. Let \((A_*, d_A)\) be a chain complex. Show if there exists a chain homomorphism \( f_* : (A_*, d_*) \to (A_*, d_*) \) chain homotopic to the identity on \((A_*, d_*)\), then \((A_*, d_*)\) is acyclic.

**Hom.** Exercises on \( \text{Hom}_R(\cdot, \cdot) \).

21. Let \( R \) be a commutative ring and \( M, N \) \( R \)-modules. Recall that \( \text{Hom}_R(M, N) \) is an \( R \)-module by Example 37.6(9). If \( h : A \to B \) is an \( R \)-homomorphism of \( R \)-modules, define
   
   \[ h_* : \text{Hom}_R(N, A) \to \text{Hom}_R(N, B) \text{ by } f \mapsto h \circ f \text{ and } h^* : \text{Hom}_R(B, N) \to \text{Hom}_R(A, N) \text{ by } f \mapsto f \circ h. \]

   Show that these are \( R \)-homomorphisms and if
   
   \[ 0 \to A \xrightarrow{f} B \xrightarrow{g} C \to 0 \]

   is a short exact sequence of \( R \)-modules and \( R \)-homomorphisms, then
   
   \[ 0 \to \text{Hom}_R(N, A) \xrightarrow{f_*} \text{Hom}_R(N, B) \xrightarrow{g_*} \text{Hom}_R(N, C) \text{ and } \]

   \[ 0 \to \text{Hom}_R(C, N) \xrightarrow{g^*} \text{Hom}_R(B, N) \xrightarrow{f^*} \text{Hom}_R(A, N) \]

   are exact. (Note the missing 0 on the right.) If \( R \) is not commutative, then your proof will only show that \( h_* \) and \( h^* \) are abelian group homomorphisms and the sequences are exact sequences of abelian groups.

22. Let \( R \) be a commutative ring and \( M_i, i \in I, \) and \( N \) be \( R \)-modules. Using the Universal Property of Direct Sums (cf. Exercise 37.19(9)), show if \( M_i, i \in I, \) and \( N \) are \( R \)-modules, then there exists an \( R \)-isomorphism \( \text{Hom}_R(\coprod_I M_i, N) \cong \prod_I \text{Hom}_R(M_i, N) \). If \( R \) is not commutative, then the isomorphism is only an abelian group isomorphism.

23. Let \( R \) be a commutative ring and \( M \) and \( N_i, i \in I \) be \( R \)-modules. Using the Universal Property of Direct Products (cf. Exercise 37.19(10)), show if \( M_i, i \in I \) and \( N \) are \( R \)-modules, then there exists an \( R \)-isomorphism \( \text{Hom}_R(M_i, \prod_I N_i) \cong \prod_I \text{Hom}_R(M_i, N) \). Remember that if \( I \) is finite, then the coproduct and product are equal. If \( R \) is not commutative, then the isomorphism is only an abelian group isomorphism.
24. Let $R$ be a (commutative) ring, $M$ and $N$ $R$-modules. Then $\text{Hom}_R(M, \ )$ and $\text{Hom}_R(\ , N)$ take split exact sequences to split exact sequences.

**Injective Modules.** The following exercises develop the theory of injective modules.

25. Let $Q$ be an $R$-module. Then $Q$ is called $R$-injective if given any $R$-monomorphism $f : A \to B$ and $R$-homomorphism $g : A \to Q$, there exists an $R$-homomorphism $h : B \to Q$ such that the diagram

$$
\begin{array}{ccc}
A & \xrightarrow{f} & B \\
\downarrow{g} & & \downarrow{h} \\
Q & & 
\end{array}
$$

commutes. Show that $Q$ is an $R$-injective if and only if, whenever

$$0 \to A \xrightarrow{f} B \xrightarrow{g} C \to 0$$

is a short exact sequence of $R$-modules and $R$-homomorphisms, then

$$0 \to \text{Hom}_R(C, Q) \xrightarrow{g^*} \text{Hom}_R(B, Q) \xrightarrow{f^*} \text{Hom}_R(A, Q) \to 0$$

is exact. (Cf. Exercise (21).)

26. Let $Q_i, i \in I$ be $R$-modules. Show that $\prod_I Q_i$ is $R$-injective if and only if $Q_i$ is $R$-injective for all $i \in I$. In particular, if $I$ is finite, then $\coprod_I Q_i$ is $R$-injective if and only $Q_i$ is $R$-injective for all $i \in I$ (since then the coproduct and product of the $Q_i$’s are the same).

27. Let $A$ be an $R$-module. Show that $A$ is $R$-injective if and only if any short exact sequence of $R$-modules of the form $0 \to A \xrightarrow{f} B \xrightarrow{g} C \to 0$ splits.

28. (Baer Criterion) Let $Q$ be an $R$-module. Show that $Q$ is an $R$-injective if and only if given any ideal $\mathfrak{A}$ in $R$ and an $R$-homomorphism $g : \mathfrak{A} \to Q$, there exists an $R$-homomorphism $h : R \to Q$ such that the diagram

$$
\begin{array}{ccc}
\mathfrak{A} & \xrightarrow{\text{inc}} & R \\
\downarrow{g} & & \downarrow{h} \\
Q & & 
\end{array}
$$

commutes where inc is the inclusion.

29. Let $R$ be a domain with $F$ its quotient field. Use the Baer Criterion to show that $F$ is $R$-injective.

30. Show that a divisible abelian group is an injective $\mathbb{Z}$-module.
31. Let $R$ be a PID. Show that an $R$-module $M$ is $R$-injective if and only if the map $\lambda_r : M \to M$ defined by $m \mapsto rm$ is an $R$-epimorphism for all nonzero $r \in R$.

32. Let $R$ be a ring and $D$ a divisible abelian group. Show that $\text{Hom}_Z(R, D)$ is an $R$-module and is $R$-injective.

33. Let $R$ be a ring and $M$ an injective $\mathbb{Z}$-module. Show that $\text{Hom}_Z(R, D)$ is an injective $R$-module.

34. Let $R$ be a ring and $M$ an $R$-module. Show that there exists an injective $R$-module $N$ and an $R$-monomorphism $\phi : M \to N$.

35. A cochain complex $(A_\ast, d)$ is called positive if $A_n = 0$ for all $n < 0$. Let $R$ be a ring and $M$ an $R$-module. Let $(A_\ast, d)$ be a positive acyclic cochain complex of $R$-modules. If we have a cochain complex $0 \to M \to A_0 \to A_1 \to \cdots$, we call this an acyclic resolution of $M$ with augmentation $\epsilon$. We write this as $0 \to M \to A_\ast$ is an acyclic resolution of $M$. If $(A_\ast, d)$ consists of injective $R$-modules, then we call such an acyclic resolution an injective resolution of $M$. Let $M$ be an $R$-module. Show that an injective resolution of $M$ exists. Show this is constructed so that $0 \to \text{coker} \, \epsilon \to A_1 \to \cdots$ is an injective resolution of $\text{coker} \, \epsilon$.

36. (Comparison Theorem) Suppose we have a diagram of positive exact complexes

$$
\begin{array}{ccl}
0 & \longrightarrow & M \\
\downarrow f & & \downarrow f^* \\
0 & \longrightarrow & X_\ast \\
\end{array}
$$

with $(I_\ast, d_\ast)$ an injective resolution of $M$. Prove there exists a cochain homomorphism $f^* : (X_\ast, d_X^\ast) \to (I_\ast, d_I^\ast)$ that satisfies the following:

(i) The diagram

$$
\begin{array}{ccl}
0 & \longrightarrow & M \\
\downarrow f & & \downarrow f^* \\
0 & \longrightarrow & X_\ast \\
\end{array}
$$

commutes.

(ii) The cochain map $f^*$ is unique up to chain homotopy.

37. Let $M$ be an $R$-module and $0 \to M \to I_\ast$ an $R$-injective resolution. [If $M = 0$, then $0 \to M \to 0$ is an injective resolution of $M$.] If $N$ is an $R$-module, let $(\text{Hom}_R(I_\ast, N), d_{\text{Hom}_R(I_\ast, N)})$ be the induced
cochain complex of abelian groups. This is a cochain complex of \( R \)-modules if \( R \) is commutative. Define the abelian groups \( (R\text{-modules if } R \text{ is commutative}) \) \( H^n(I_M, N) \), \( n \geq 0 \), to be the cohomology of this cochain complex. Prove that \( H^n(I_M, N) \) is independent of the injective resolution \( I_{M,*} \) of \( M \). [Actually, any two injective resolutions give cohomology groups up to a “natural” unique isomorphism between them as your proof will show, so we identify them.] We denote the cohomology group \( H^n(I_M, N) \) by \( \text{Ext}^n_R(M, N) \). Note that \( \text{Ext}^0_R(M, N) = \text{Hom}_R(M, N) \).

38. (Horseshoe Lemma) Let

\[
\begin{array}{ccccccc}
0 & \longrightarrow & M' & \overset{f}{\longrightarrow} & M & \overset{g}{\longrightarrow} & M'' & \longrightarrow & 0 \\
\end{array}
\]

be a short exact sequence of \( R \)-modules and

\[
\begin{array}{ccccccc}
0 & \longrightarrow & M' & \overset{\varepsilon'}{\longrightarrow} & I_{*} & \overset{\varepsilon}{\longrightarrow} & I_{*}'' & \longrightarrow & 0 \\
\end{array}
\]

injective resolutions. Then there exist an injective resolution

\[
\begin{array}{ccccccc}
0 & \longrightarrow & M & \overset{\varepsilon}{\longrightarrow} & I_{*} & \\
\end{array}
\]

such that

\[
\begin{array}{ccccccc}
0 & \longrightarrow & I_{*} & \overset{f^*}{\longrightarrow} & I_{*} & \overset{g^*}{\longrightarrow} & I_{*}'' & \longrightarrow & 0 \\
\end{array}
\]

is a split exact sequence of injective chain complexes, i.e.,

\[
\begin{array}{ccccccc}
0 & 0 & 0 & 0 \\
\downarrow & \downarrow & \downarrow & \downarrow \\
0 & \longrightarrow & M' & \overset{f}{\longrightarrow} & M & \overset{g}{\longrightarrow} & M'' & \longrightarrow & 0 \\
\varepsilon' & \downarrow & \varepsilon & \downarrow & \varepsilon'' \\
0 & \longrightarrow & I_{0}' & \overset{f^0}{\longrightarrow} & I_{0} & \overset{g^0}{\longrightarrow} & I_{0}'' & \longrightarrow & 0 \\
\varepsilon_{0}' & \downarrow & \varepsilon_{0} & \downarrow & \varepsilon_{0}'' \\
0 & \longrightarrow & I_{1}' & \overset{f^1}{\longrightarrow} & I_{1} & \overset{g^1}{\longrightarrow} & I_{1}'' & \longrightarrow & 0 \\
\vdots & \vdots & \vdots & \vdots \\
\end{array}
\]

commutes and has exact columns and split exact rows (except for the top row).

39. Let \( N \) be an \( R \)-module. If

\[
\begin{array}{ccccccc}
0 & \rightarrow & A & \overset{f}{\rightarrow} & B & \overset{g}{\rightarrow} & C & \rightarrow & 0 \\
\end{array}
\]
is a short exact sequence of $R$-modules and $R$-homomorphisms, show that the exact sequence

$$0 \to \text{Hom}_R(A, N) \xrightarrow{f} \text{Hom}_R(B, N) \xrightarrow{g} \text{Hom}_R(C, N)$$

extends to a long exact sequence in cohomology

$$\cdots \to \text{Ext}^{n-1}_R(C, N) \xrightarrow{\partial^{n-1}} \text{Ext}^n_R(A, N) \xrightarrow{\tau^n} \text{Ext}^n_R(B, N) \xrightarrow{\tau^n} \text{Ext}^n_R(C, N) \xrightarrow{\partial^n} \text{Ext}^{n+1}_R(A, N) \to \cdots$$

40. Let $M$ and $N$ be $R$-modules. If $N$ is $R$-injective, show $\text{Ext}^n_R(M, N) = 0$ for all $n > 0$.

41. Let $A$ and $M$ be $R$-modules. Suppose that

$$0 \to M \xrightarrow{f} I \xrightarrow{g} N \to 0$$

and

$$0 \to M' \xrightarrow{f'} I' \xrightarrow{g'} N' \to 0$$

are short exact sequences of $R$-modules with $I$ and $I'$ injective $R$-modules. Show the following:

(i) (Schanuel’s Lemma) $N \oplus I' \cong N' \oplus I$.

(ii) (Dimension Shift) $\text{Ext}^{n+1}_R(A, M) = \text{Ext}^n_R(A, N) = \text{Ext}^n_R(A, N')$ for all $n > 1$.

42. Let $M, N_i$ be $R$-modules with $i \in I$. Show that

$$\text{Ext}^*_R(M, \prod_I N_i) = \prod_I \text{Ext}^*_R(M, N_i).$$

Note if $I$ is finite, then $\text{Ext}^*_R(M, \coprod_I N_i) = \coprod_I \text{Ext}^*_R(M, N_i)$.

43. Let $N$ and $R$-module. Show the following are equivalent:

(i) $N$ is $R$-injective.

(ii) $\text{Ext}^n_R(M, N) = 0$ for all all $n \geq 1$ and all $R$-modules $M$.

(iii) $\text{Ext}^1_R(M, N) = 0$ for all $R$-modules $M$.

38. Free Modules

Vector spaces have nice properties because of the concept of linear independence and the fact that one can divide by any nonzero scalar. In general, $R$-modules have neither of these two properties. There is little we can do about the lack of units in an arbitrary ring, but we can look at $R$-modules that are spanned by linearly independent sets. Such modules are called free $R$-modules. Though sparse, these modules are quite nice, in fact, just as a vector space is isomorphic to a external direct sum of copies of the underlying field, the analogous statement is true for free $R$-modules. The first result one would then expect is that the notion of dimension makes sense, i.e., any two linearly independent
spanning set of a free module should have the same cardinality. Unfortunately, this is false in general. The problem does vanish, however, if the ring $R$ is commutative. We shall leave the proof of this fact as an exercise. This is fortuitous as these are the rings we are studying in detail.

**Definition 38.1.** A nonzero $R$-module $M$ is called a free $R$-module (or $R$-free) if there exists a basis for $M$, i.e., a subset $\mathcal{B}$ of $M$ satisfying:

(i) $M = \langle \mathcal{B} \rangle$, i.e., $\mathcal{B}$ generates or spans $M$.

(ii) $\mathcal{B}$ is linearly independent, i.e., if $0 = \sum_{x \in \mathcal{B}} r_x x$ (so $r_x = 0$ for almost all $x \in \mathcal{B}$), then $r_x = 0$ for all $x \in \mathcal{B}$. Equivalently, $M = \bigoplus_{x \in \mathcal{B}} Rx$ and if $rx = 0$ with $x \in \mathcal{B}$, $r \in R$, then $r = 0$.

We shall also call the trivial module a free $R$-module.

Of course, a set $\mathcal{D}$ in $M$ is called linearly dependent if it is not linearly independent. Equivalently, there exist $x_1, \ldots, x_n \in \mathcal{D}$, some $n$, and $r_1, \ldots, r_n$ in $R$ not all zero such that $r_1 x_1 + \cdots + r_n x_n = 0$.

If $M$ is a free $R$-module on basis $\mathcal{B}$, the definition says that $\mathcal{B}$ generates $M$ and there are no nontrivial relations on these generators. Contrast this with a finite (abelian) group say of order $n$ which says that every element $x$ satisfies the relation $x^n = e$. In particular, this is true of any element in any generating set.

**Examples 38.2.** Let $R$ be a ring and $M$ be a nontrivial $R$-module.

1. $R$ is $R$-free on basis $\{u\}$ with $u$ a unit in $R$.

2. Suppose that $M$ is $R$-cyclic. Then $M$ is $R$-free if there exists an $x \in M$ satisfying $M = Rx$ and $\text{ann}_R x = 0$. It follows that $M$ is $R$-free if and only if $M \cong R$.

3. If $M$ is a finitely generated $R$-module, then $M$ is $R$-free if and only if there exists a finite basis for $M$ if and only if

$$M \cong R \bigoplus \cdots \bigoplus R = R^n$$

for some positive integer $n$. (Why?)

4. If $R$ is a field (or a division ring), then every $R$-module is free and any two bases have the same cardinality.

5. If $R$ is commutative and $M$ is $R$-cyclic, then $M$ is $R$-free if and only if $\text{ann}_R m = 0$ whenever $M = Rm$, i.e., for any generator of $M$.

To show this one needs to use the first two parts of Exercise 38.12(7) (which establishes that all bases for a finitely generated free module over a commutative ring have the same cardinality). If we know this and $M$ is $R$-free, then say every basis for $M$ has say $n$ elements. Let
\[ B \text{ be a basis for } M \text{ and } m \text{ be a maximal ideal in } R. \] Then this exercise would allow us to prove the following sequence of isomorphisms:

\[ Rm/mm = M/mm = \left( \bigoplus_m Rx \right)/\left( m\left( \bigoplus_m Rx \right) \right) \]

\[ \cong R^n/mR^n \cong (R/m)^n \]

as \( R \)- and as \( R/m \)-modules. As \( R/m \) is a field, this is an isomorphism of vector spaces over \( R/m \). As \( m + mm \) generates \( Rm/mm \), by (3), we must have \( n = 1 \).

6. If \( R \) is not commutative, it can happen that \( R^n \cong R^m \) for different positive integers, i.e., different bases can have different cardinalities. For example, let \( M \) be a free \( R \)-module on a countable (not finite) basis (e.g., \( M = R[t] \) (cf (9) below)). Let \( S \) be the endomorphism ring \( \text{End}_R(M) \). Then it can be shown that \( S^n \) is \( S \)-free and \( S^n \cong S^m \) for all positive integers \( n \) and \( m \).

7. Suppose that \( M \) is \( R \)-free with basis \( B \) and \( x \in B \). Then the map \( \rho_x : R \rightarrow M \) given by \( r \mapsto rx \) is an \( R \)-monomorphism. In particular, the cardinality of \( M \) is at least that of \( R \). This gives another proof that no nontrivial finite abelian group is \( Z \)-free.

8. Let \( M_i \) be free \( R \)-modules for \( i \in I \). Then \( \bigoplus_I M_i \) is \( R \)-free. [If \( B_i \) is a basis for \( M_i \) for each \( i \in I \), what is a basis for \( \bigoplus_I M_i \)?] For example \( \bigoplus_{i=1}^{\infty} Z \) is \( Z \)-free on the standard basis

\[ S = \{ e_i = (0, \ldots, 1, 0, \ldots) \mid i \geq 1 \}. \]

[Warning: It turns out that \( \bigoplus_{i=0}^{\infty} Z \) is not \( Z \)-free.]

9. \( Q \) is not \( Z \)-free.

10. \( R[t] \) is \( R \)-free on basis \( B = \{ t^i \mid i \geq 0 \} \). One can now define the ring \( R[t] \) to be the free \( R \)-module on basis \( B \) with multiplication induced by \( t^i \cdot t^j := t^{i+j} \) (and extend linearly). More generally, \( R[t_1, \ldots, t_n] \) is the free \( R \)-module on basis \( \{ t_1^{i_1} \cdots t_n^{i_n} \mid i_j \geq 0 \text{ for all } i_j \} \).

11. Let \( R_0 \) be a ring and \( R = R_0[t_1, t_2] \). The \( R \)-module \( M = Rt_1 + Rt_2 \) is a submodule of \( R \) as it is a left ideal. Then \( M \) is not \( R \)-free. [Proof?]

[This shows that a submodule of a free \( R \)-module may not be free.]

12. Let \( n \mid m \in \mathbb{Z}^+ \) and \( - : \mathbb{Z}/m\mathbb{Z} \rightarrow \mathbb{Z}/n\mathbb{Z} \) be the ring epimorphism given by \( r + m\mathbb{Z} \mapsto r + n\mathbb{Z} \). Then \( \mathbb{Z}/n\mathbb{Z} \) is a \( (\mathbb{Z}/m\mathbb{Z}) \)-module via the pullback, e.g., \( \mathbb{Z}/2\mathbb{Z} \) and \( \mathbb{Z}/3\mathbb{Z} \) are \( (\mathbb{Z}/6\mathbb{Z}) \)-modules. By element count neither \( \mathbb{Z}/2\mathbb{Z} \) nor \( \mathbb{Z}/3\mathbb{Z} \) are \( \mathbb{Z}/6\mathbb{Z} \)-free, but \( \mathbb{Z}/2\mathbb{Z} \uplus \mathbb{Z}/3\mathbb{Z} \cong \mathbb{Z}/6\mathbb{Z} \) as \( (\mathbb{Z}/6\mathbb{Z}) \)-modules (why?), so a direct summand of a free
module may not be free. (A submodule $N$ is a *direct summand* of an $R$-module $M$ if $M = N \oplus M'$ for some submodule $M'$.)

13. If $M = \bigoplus_{i} R x_i$ is a free $R$-module on basis $B = \{x_i \mid i \in I\}$ and $m$ is an element in $M$, then there exist unique $r_i \in R, i \in I$, almost all zero such that $m = \sum_{i} r_i x_i$. The $r_i$ are called the *coordinates* of $m$ relative to the basis $B$ and $r_i$ is called the *coordinate* of $m$ on $x_i$.

14. The set $\{2\}$ is linearly independent in the free $\mathbb{Z}$-module $\mathbb{Z}$, but it is not a basis as $2\mathbb{Z} \subset \mathbb{Z}$. Indeed any two distinct elements in $\mathbb{Z}$ form a linearly dependent set. This shows that, in general, not every linearly independent set in a free $R$-module need be part of a basis.

15. Let $R$ be a domain and $M$ a free $R$-module. If $z$ is a nonzero element in $M$ satisfying $rz = 0$ for some $r \in R$, then $r = 0$ (for if $B$ is a basis, then $z = \sum_{i} r x_i x_i$).

Universal properties usually produce maps satisfying certain properties together with a uniqueness statement. This is very useful, as one of the more difficult problems is to show the existence of a map, or to prove that a potential explicit map is well-defined. Usually the problem arises as we know generators but cannot show that the putative map respects the relations among these generators. As a basis for a free module does not satisfy any nontrivial relations, this causes no problems and free modules satisfy the following universal property:

**Theorem 38.3.** (Universal Property of Free Modules) Let $B = \{x_i\}_I$ be a basis for a free $R$-module $M$. If $N$ is an $R$-module and $y_i, i \in I$, elements in $N$ (not necessarily distinct), then there exists a unique $R$-homomorphism $f : M \to N$ such that $x_i \mapsto y_i$ for all $i \in I$.

**Proof.** If $z \in M$, there exist unique $r_i \in R$, almost all $r_i = 0$, such that $z = \sum_{i} r_i x_i$. In particular, the uniqueness of the $r_i$ implies that $f : M \to N$ given by $z \mapsto \sum_{i} r_i y_i$ is well-defined. Clearly, $f$ is uniquely determined by $x_i \mapsto y_i$ and $f$ is an $R$-homomorphism.

In the notation of the theorem, we say that $x_i \mapsto y_i$ extends linearly to a homomorphism $f : M \to N$. The theorem says that any $R$-homomorphism from a free module $M$ to another $R$-module is completely determined by where a basis for $M$ is sent. Note, of course, that the variant of this statement (and proof) is one that you should know about vector spaces, which are just free modules over a field.

**Remark 38.4.** A better way of writing the universal property of free modules is the following: Let $M$ be a free $R$-module on basis $B$ and $N$ an $R$-module. Given any set map $g : B \to N$ there exists a unique
$R$-homomorphism $f : M \to N$ such that the diagram

$$
\begin{array}{ccc}
\mathcal{B} & \xrightarrow{\text{inc}} & M \\
\downarrow{g} & & \downarrow{f} \\
& N.
\end{array}
$$

commutes where $\text{inc}$ is the inclusion map.

**Corollary 38.5.** Let $M$ and $N$ be free $R$-modules on bases $\mathcal{B}$ and $\mathcal{C}$ respectively. If there exists a bijection $g : \mathcal{B} \to \mathcal{C}$, i.e., $|\mathcal{B}| = |\mathcal{C}|$, then $M \cong N$.

**Proof.** The maps $g$ and $g^{-1}$ of sets induce inverse $R$-isomorphisms $M \to N$ and $N \to M$. 

**Corollary 38.6.** Let $M$ be an $R$-module (respectively, a finitely generated $R$-module). Then there exists a free $R$-module (respectively, a finitely generated free $R$-module) $P$ and an $R$-epimorphism $g : P \to M$. In particular, we have a short exact sequence

$$
0 \to \ker g \to P \stackrel{g}{\to} M \to 0.
$$

**Proof.** Let $Y = \{y_i\}_I$ generate $M$. If $M$ is finitely generated, we may assume that $I$ is finite. Let $P = \bigsqcup_I R$ and $S_I$ the standard basis for $P$. (So $S_I := \{e_i \mid e_i = (\delta_{ij})_{j \in I}\}_I$, i.e., $e_i$ has 1 in the $i$th component, 0 elsewhere.) Then the set map $S_I \to M$ given by $e_i \mapsto y_i$ extends linearly to an $R$-homomorphism $g : P \to M$. As $Y$ generates $M$, the map is surjective. Unfortunately, the kernel of $g$ in the corollary may not be nice, in particular, it may not be free. If $R$ is a PID, it will be free. We shall show this later under the assumption that $M$ is finitely generated.

**Example 38.7.** Let $R$ be a commutative ring and $M$ a free $R$-module on basis $\mathcal{B}$. As $R$ is commutative, $M^* := \text{Hom}_R(M, R)$ is an $R$-module, called the dual module of $M$. [Warning: It is usually not $R$-free.] For each $x \in \mathcal{B}$ define the $x$th coordinate function $f_x : M \to R$ by $\sum_R r_x x \mapsto r_x$. It is an $R$-homomorphism. Set $\mathcal{B}^* := \{f_x \mid x \in \mathcal{B}\}$. The map $\mathcal{B} \mapsto \mathcal{B}^*$ by $x \mapsto f_x$ extends linearly to an $R$-homomorphism $f : M \mapsto M^*$. This map is monic and its image is $R$-free on basis $\mathcal{B}^*$. We call $\mathcal{B}^*$ the dual basis of $\mathcal{B}$. Unfortunately, this map depends on the basis $\mathcal{B}$, so it is not “natural”. However, if we let $M^{**} := \text{Hom}_R(M^*, R)$ and $L : M \to M^{**}$ the evaluation map at each $x$ in $M$, i.e., the map $L$ is given by $x \mapsto L_x : f \mapsto f(x)$, then this map is an $R$-monomorphism and is “natural” as it does not depend on the choice of a basis. If $\mathcal{B}$ is finite, then $f$ is also surjective, hence $M \cong M^*$ (and $M \cong M^{**}$).
We end this section with results left as important exercises.

**Definition 38.8.** We say that a ring \( R \) satisfies the *invariant dimension property* or IDP if for every finitely generated free \( R \)-module \( P \), every basis for \( P \) has the same cardinality (finite by Example 38.2(3)). If \( R \) satisfies IDP, and \( P \) is a finitely generated free \( R \)-module, then the cardinality of a basis for \( P \) is called the *rank* of \( P \) and written \( \text{rank} \ P \).

We know that fields satisfy IDP. In fact, we have shown this (cf. Remark 17.13). It is a main ingredient in the proof of the following theorem that we leave as an exercise (cf. Exercise 38.12(7)).

**Theorem 38.9.** Let \( R \) be a commutative ring. Then \( R \) satisfies IDP, i.e., if \( M \) is a finitely generated free \( R \)-module then all bases of \( M \) have the same cardinality.

We also leave the proof of the following corollary to this theorem as an exercise (Exercise 38.12(8)).

**Corollary 38.10.** Let \( R \) be a commutative ring and \( M \) and \( N \) be finitely generated free \( R \)-modules. Then \( \text{rank} \ M \oplus N = \text{rank} \ M + \text{rank} \ N \).

**Remark 38.11.** It can be shown for any ring \( R \), if \( M \) is a free \( R \)-module on basis \( B \) with \( B \) not finite, i.e., \( M \) is not finitely generated, then every basis of \( M \) has the same cardinality.

In the exercises, we continue to state some basic results in the study of (co)homology. These exercises may be safely omitted until the reader is more comfortable with module theory.

**Exercises 38.12.**

1. Show if \( M \) is a nonzero finitely generated \( R \)-module, then \( M \) is \( R \)-free if and only if there exists a finite basis for \( M \) if and only if \( M \cong R^n \) for some positive integer \( n \).

2. Show that \( \mathbb{Q} \) is not \( \mathbb{Z} \)-free. (Hint: Show any two distinct elements of \( \mathbb{Q} \) are \( \mathbb{Z} \)-linearly dependent.)

3. Let \( R_0 \) be a commutative ring and \( R = R_0[t_1, t_2] \). Let \( M = Rt_1 + Rt_2 \) Show the \( R \)-module \( M \) is not \( R \)-free.

4. Let \( R \) be a nontrivial commutative ring. Show that \( R \) is a field if and only if every finitely generated \( R \)-module is free.

5. Suppose that

\[
0 \to A \xrightarrow{f} B \xrightarrow{g} C \to 0
\]

is a short exact sequence of \( R \)-modules with \( C \) a free \( R \)-module. Show the sequence splits. (Cf. Exercise 37.19(12) for the definition of a split exact sequence.)
6. Let $M$ be a free $R$-module. Suppose that $\mathcal{B}$ is a set and $\iota : \mathcal{B} \rightarrow M$ is set map (not necessarily one-to-one). Then $\iota : \mathcal{B} \rightarrow M$ satisfies the following universal property: Given any set map $j : \mathcal{B} \rightarrow N$, there exists a unique $R$-homomorphism $f : M \rightarrow N$ such that the diagram

\[
\begin{array}{ccc}
\mathcal{B} & \xrightarrow{\iota} & M \\
\downarrow g & & \downarrow j \\
N.
\end{array}
\]

commutes. [Cf. Remark 38.4.]

7. Let $P$ be a free $R$-module on basis $\mathcal{B} = \{x_i\}_{i \in I}$. Let $\mathfrak{A} < R$ be a (2-sided) ideal. Show all of the following:

(i) $P/\mathfrak{A}P \cong \prod_I Rx_i/\mathfrak{A}x_i \cong \prod_I R/\mathfrak{A}$.

(ii) Let $\overline{\cdot} : R \rightarrow R/\mathfrak{A}$ be the canonical ring epimorphism. Set $\overline{\mathcal{B}} = \{\overline{x_i} := x_i + \mathfrak{A}P \mid i \in I\}$. Then $P/\mathfrak{A}P$ is a free $\overline{R}$-module on basis $\overline{\mathcal{B}}$ and $|\overline{\mathcal{B}}| = |\mathcal{B}|$.

(iii) Let $\varphi : R \rightarrow S$ be a ring epimorphism with $S \neq 0$. If $S$ satisfies IDP so does $R$.

(iv) Any commutative ring satisfies IDP.

8. Prove Corollary 38.10.

9. Show Example 38.2(6) has the desired properties.

[Hint: If $\mathcal{B} = \{e_i \mid i \in \mathbb{Z}^+\}$ is a basis for $P$, show that $\{f_1, f_2\}$ is a basis for $S$, where for all $n$, we have $f_1(e_{2n}) = e_n$, $f_1(e_{2n+1}) = 0$ and $f_2(e_{2n}) = 0$, $f_2(e_{2n+1}) = e_n$.]

**Projective Modules.** The following exercises develop the theory of projective modules.

10. Let $P$ be an $R$-module. Then $P$ is called $R$-projective if given any $R$-epimorphism $f : B \rightarrow C$ and $R$-homomorphism $g : P \rightarrow C$, there exists an $R$-homomorphism $h : P \rightarrow B$ such that the diagram

\[
\begin{array}{ccc}
P & \xrightarrow{h} & B \\
\downarrow g & & \downarrow f \\
C.
\end{array}
\]

commutes. Show that any free $R$-module is projective.

11. Show that a direct summand of an $R$-free module is projective and a direct sum of $R$-modules is projective if and only if each direct summand of it is $R$-projective.
12. Show \( \mathbb{Z}/2\mathbb{Z} \) is a projective \( \mathbb{Z}/6\mathbb{Z} \)-module but not is not \( \mathbb{Z}/6\mathbb{Z} \)-free.

13. Let \( P \) be an \( R \)-module. Show that \( P \) is a projective \( R \)-module if and only if, whenever

\[
(*) \quad 0 \rightarrow A \xrightarrow{f} B \xrightarrow{g} C \rightarrow 0
\]

is a short exact sequence of \( R \)-modules and \( R \)-homomorphisms, then

\[
0 \rightarrow \text{Hom}_R(P, A) \xrightarrow{f^*} \text{Hom}_R(P, B) \xrightarrow{g^*} \text{Hom}_R(P, C) \rightarrow 0
\]

is exact. In particular, if \( C \) is \( R \)-projective, then (*) is split exact. (Cf. Exercises 37.19(21) and 37.19(25)).

14. Let \( C \) be an \( R \)-module. Show that \( C \) is \( R \)-projective if and only if any short exact sequence of \( R \)-modules of the form

\[
0 \rightarrow A \xrightarrow{f} B \xrightarrow{g} C \rightarrow 0
\]

splits.

15. (Projective Basis) Let \( M \) be an \( R \)-module. Prove that \( M \) is \( R \)-projective if and only if there exist a set of elements \( \{m_i \mid m_i \in M \text{ and } i \in I\} \) and a set of \( R \)-homomorphisms \( \{\varphi_i : A \rightarrow R \mid i \in I\} \) such that for all \( m \in M \), we have \( m = \sum_I \varphi_i(m)m_i \) with almost all \( \varphi(m) = 0 \). Moreover, if this is the case, then \( \{m_i \mid i \in I\} \) generate \( M \). [Cf. this a basis for a free module.]

16. A chain complex \( (A_*, d_*) \) is called positive if \( A_n = 0 \) for all \( n < 0 \). Let \( R \) be a ring and \( M \) an \( R \)-module. Let \( (A_*, d_*) \) be a positive acyclic chain complex of \( R \)-modules. If we have a chain complex

\[
\cdots \xrightarrow{d_2} A_1 \xrightarrow{d_1} A_0 \xrightarrow{\varepsilon} M \rightarrow 0
\]

we call this an acyclic resolution of \( M \) with augmentation \( \varepsilon \). We write this as \( A_* \xrightarrow{\varepsilon} M \rightarrow 0 \) is an acyclic resolution of \( M \). If \( (A_*, d_*) \) consists of projective \( R \)-modules, then we call such an acyclic resolution a projective resolution of \( M \). Let \( M \) be an \( R \)-module. Show that a projective resolution of \( M \) exists. Show this is constructed so that \( \cdots \xrightarrow{(d_A)_2} A_1 \xrightarrow{(d_A)_1} \ker \varepsilon \rightarrow 0 \) is a projective resolution of \( \ker \varepsilon \).

17. (Comparison Theorem) Suppose we have a diagram of positive exact complexes

\[
\begin{array}{ccc}
P_* & \xrightarrow{\varepsilon} & M \xrightarrow{f} 0 \\
\downarrow f & & \downarrow \\
X_* & \xrightarrow{\varphi} & M \xrightarrow{g} 0
\end{array}
\]

with \( (P_*, d_{P*}) \) a projective resolution of \( M \). Show that there exists a chain homomorphism \( f_* : (P_*, d_{P*}) \rightarrow (X_*, d_{P*}) \) that satisfies the following:
(i) The diagram

\[
\begin{array}{ccc}
P_* & \xrightarrow{\varepsilon} & M \\
\downarrow{f_*} & & \downarrow{f} \\
X_* & \xrightarrow{\varphi} & M
\end{array}
\]

commutes.

(ii) The chain map \( f_* \) is unique up to chain homotopy.

18. Let \( N \) be an \( R \)-module and \( P_{N_*} \xrightarrow{\varepsilon} N \to 0 \) an \( R \)-projective resolution. [If \( N = 0 \), then \( 0 \to N \xrightarrow{0} 0 \) is a projective resolution of \( N \).] If \( M \) is an \( R \)-module, let \( (\text{Hom}_R(M, P_{N_*}), d_{\text{Hom}_R(M, N_P)}) \) be the induced cochain complex of abelian groups. This is a cochain complex of \( R \)-modules if \( R \) is commutative. Define the abelian groups (\( R \)-modules if \( R \) is commutative) \( H^n(M, P_N) \), \( n \geq 0 \), to be the cohomology of this cochain complex. Prove that \( H_n(M, P_N) \) is independent of the projective resolution of \( N \). [Actually, any two projective resolutions give cohomology groups up to a “natural” unique isomorphism between them as your proof will show, so we identify them.) We denote \( H_*(M, N) \) by \( \text{Ext}_R^*(M, N) \). Note that \( \text{Ext}_R^0(M, N) = \text{Hom}_R(M, N) \).

19. (Horseshoe Lemma) Let

\[
\begin{array}{ccc}
0 & \longrightarrow & M' \\
& \xrightarrow{f} & M \\
& & \xrightarrow{g} M'' \\
& & \longrightarrow 0
\end{array}
\]

be a short exact sequence of \( R \)-modules and

\[
\begin{array}{cc}
P_*' & \xrightarrow{\varepsilon'} M' \\
& \longrightarrow 0 \\
& \quad \text{and} \\
P_*'' & \xrightarrow{\varepsilon'} M'' \\
& \longrightarrow 0
\end{array}
\]

projective resolutions. Then there exist an \( R \)-projective resolution

\[
P_* \xrightarrow{\varepsilon} M \longrightarrow 0
\]

such that

\[
\begin{array}{ccc}
0 & \longrightarrow & P_*' \\
& \xrightarrow{f_*} & P_* \\
& & \xrightarrow{g_*} P_*'' \\
& & \longrightarrow 0
\end{array}
\]

is a split exact sequence of projective chain complexes, i.e.,
20. Let $M$ be an $R$-module. If

$$0 \rightarrow A \xrightarrow{f} B \xrightarrow{g} C \rightarrow 0$$

is a short exact sequence of $R$-modules and $R$-homomorphisms, show that the exact sequence of abelian groups (respectively, $R$-homomorphisms if $R$ is commutative)

$$0 \rightarrow \text{Hom}_R(M, C) \xrightarrow{g^*} \text{Hom}_R(M, B) \xrightarrow{f^*} \text{Hom}_R(M, A)$$

extends to a long exact sequence of abelian groups (respectively $R$-modules if $R$ is commutative) in cohomology

$$\cdots \rightarrow \text{Ext}^{n-1}_R(M, A) \xrightarrow{\partial^{n-1}} \text{Ext}^n_R(M, C) \xrightarrow{\partial^n} \text{Ext}^n_R(M, B) \xrightarrow{f^n} \text{Ext}^n_R(M, A) \xrightarrow{g^n} \text{Ext}^{n+1}_R(M, C) \rightarrow \cdots$$

21. Let $M$ and $N$ be $R$-modules. If $M$ is $R$-projective, show $\text{Ext}^n_R(M, N) = 0$ for all $n > 0$.

22. Let $A$ and $M$ be $R$-modules. Suppose that

$$0 \rightarrow N \xrightarrow{f} P \xrightarrow{g} M \rightarrow 0 \text{ and } 0 \rightarrow N' \xrightarrow{f'} P' \xrightarrow{g'} M \rightarrow 0$$

are short exact sequences of $R$-modules with $P$ and $P'$ projective $R$-modules. Show the following:

(i) (Schanuel’s Lemma) $M \oplus P' \cong M' \oplus P$.

(ii) (Dimension Shift) $\text{Ext}^{n+1}_R(M, A) = \text{Ext}^n_R(N, A) = \text{Ext}^n_R(N', A)$ for all $n > 1$. 
(iii) If $M_i, N$ are $R$-modules with $i \in I$, then $\text{Ext}^*_R(\prod_I M_i, N) = \prod_I \text{Ext}^*_R(M_i, N)$. Therefore, if $I$ is finite, then $\text{Ext}^*_R(\prod_I M_i, N) = \prod_I \text{Ext}^*_R(M_i, N)$.

23. Let $M$ and $N$ be $R$-modules. Show $\text{Ext}^n_R(M, A) = \text{Ext}^n_R(M, A)$ for all $n$, so we write them both as $\text{Ext}^n_R(A, B)$. In particular, $\text{Ext}^*_R(\ , \ )$ commutes with finite direct sums in both variables.

24. Let $M$ and $R$-module. Show the following are equivalent:
   (i) $M$ is $R$-projective.
   (ii) $\text{Ext}^n_R(M, N) = 0$ for all all $n \geq 1$ and all $R$-modules $N$.
   (iii) $\text{Ext}^1_R(M, N) = 0$ for all $R$-modules $N$.

**Tensor Products** The following exercises introduce tensor products over commutative rings. [This can be generalized to arbitrary rings with suitable modifications.]

25. Let $R$ be a commutative ring and $M, N$ two $R$-modules. Let $P$ be the free $R$-module on basis $\{(m, n) \mid m \in M, n \in N\}$ and $X$ the submodule of $P$ generated by the elements
   (i) $(m_1 + m_2, n) - (m_1, n) - (m_2, n)$ for all $m_1, m_2 \in M$ and $n \in N$.
   (ii) $(m, n_1 + n_2) - (m, n_1) - (m, n_2)$ for all $m \in M$ and $n_1, n_2 \in N$.
   (iii) $(rm, n) - r(m, n)$ for all $m \in M$, $n \in N$, and $r \in R$.
   (iv) $(m, rn) - r(m, n)$ for all $m \in M$, $n \in N$, and $r \in R$.

Let $f : M \times N \to P/X$ be the $R$-bilinear map induced by $(m, n) \mapsto (m, n) + X$, i.e., an $R$ homomorphism in each variable. Show that $f : M \times N \to P/X$ satisfies the following universal property: If $g : M \times N \to Q$ is an $R$-bilinear map, then there exists a unique $R$-homomorphism $h : P/X \mapsto Q$ such that

$$
\begin{array}{ccc}
M \times N & \xrightarrow{f} & P/X \\
\downarrow{g} & & \downarrow{h} \\
Q & & 
\end{array}
$$

commutes. The $R$-module $P/X$ is called the tensor product of $M$ and $N$ and denoted by $M \otimes_R N$ and the elements $(m, n) + X$ are denoted by $m \otimes n$.

26. Let $R$ be a commutative ring and $M$ and $N$ be $R$-modules. Show that $M \otimes_R N \cong N \otimes_R M$.

27. Let $R$ be a commutative ring and $M$ an $R$-module. Show that $M \cong R \otimes_R M$.

28. Let $m$ and $n$ be positive integers with greatest common divisor $d$. Show $\mathbb{Z}/m\mathbb{Z} \otimes \mathbb{Z}/n\mathbb{Z} \cong \mathbb{Z}/d\mathbb{Z}$. 

29. Let $R$ be a commutative ring and $M_i$, $i \in I$, and $N$ be $R$-modules. Show \( (\prod_i M_i) \otimes_R N \cong \prod_i (M_i \otimes_R N) \). In particular, show if $M$ and $N$ are free $R$-modules (respectively, and finitely generated), then $M \otimes_R N$ is a free $R$-module (respectively, of rank $M \cdot \text{rank } N$).

30. Let $R$ be a commutative ring and $f_1 : M_1 \to N_1$ and $f_2 : M_2 \to N_2$ be $R$-homomorphisms. Show these induce a unique $R$-homomorphism $f_1 \otimes f_2 : M_1 \otimes_R M_2 \to N_1 \otimes_R N_2$ satisfying $m_1 \otimes m_2 \mapsto f_1(m_1) \otimes f_2(m_2)$ for all $m_1 \in M_1$, $m_2 \in M_2$.

31. Let $R$ be a commutative ring and

\[
0 \to A \xrightarrow{f} B \xrightarrow{g} C \to 0
\]

an exact sequence of $R$-modules and $R$-homomorphisms. Show for every $R$-module $M$, the sequences

\[
\begin{align*}
M \otimes_R A & \xrightarrow{1_M \otimes f} M \otimes_R B \xrightarrow{1_M \otimes g} M \otimes_R C \to 0 \\
A \otimes_R M & \xrightarrow{f \otimes 1_M} B \otimes_R M \xrightarrow{g \otimes 1_M} C \otimes_R M \to 0
\end{align*}
\]

are exact.

32. Let $R$ be a commutative ring and $A, B, C$ be $R$-modules. Show that there exists a “natural” isomorphism

\[
\theta : \text{Hom}_R(A, \text{Hom}_R(B, C)) \to \text{Hom}_R(A \otimes_R B, C),
\]

i.e., an isomorphism such that if $\varphi : A \to A'$ is an $R$-homomorphism, the there exists an $R$-homomorphism a commutative diagram

\[
\begin{array}{ccc}
\text{Hom}_R(A, \text{Hom}_R(B, C)) & \xrightarrow{\theta} & \text{Hom}_R(A \otimes_R B, C) \\
\varphi_* & & (\varphi \otimes 1_B)_* \\
\text{Hom}_R(A', \text{Hom}_R(B, C)) & \xrightarrow{\theta'} & \text{Hom}_R(A \otimes_R B, C)
\end{array}
\]

33. Let $\varphi : R \to S$ be a ring homomorphism of commutative rings. We know that $S$ becomes an $R$-module via the pull back $rs := \varphi(r)s$. Let $M$ be an $R$-module. Show that the $R$-module $S \otimes_R M$ becomes an $S$-module via $s_2(s_1 \otimes m) = s_1 s_2 \otimes m$ for all $s_1, s_2 \in S$ and $m \in M$. In particular, if $M$ is a free $R$-module on basis $B$, then $S \otimes_R M$ is a free $S$-module on basis $\{1 \otimes m \mid m \in B\}$.

34. Let $R$ be a commutative ring, $M$ an $R$-module, and $P_{M*} \xrightarrow{\partial} M \to 0$ an $R$-projective resolution of $M$. Show for every every $R$-module $N$, the homology $H^*_R(P_M \otimes_R N)$ is independent of the projective resolution of $M$. Let $\text{Tor}^n_R(M, N)$ denote $H^n_R(P_M \otimes_R N)$. Show also that if

\[
0 \to A \xrightarrow{f} B \xrightarrow{g} C \to 0
\]
is a short exact sequence of $R$-modules and $R$-homomorphisms, then the exact sequence

$$M \otimes_R A \xrightarrow{f_*} M \otimes_R B \xrightarrow{g_*} M \otimes_R A \to 0$$

extends to a long exact sequence in homology

$$\cdots \to \text{Tor}^n_{n+1}(M, C) \xrightarrow{\partial_{n+1}} \text{Tor}^n_{n+1}(M, A) \xrightarrow{f_*} \text{Tor}^n_{n+1}(M, B) \xrightarrow{g_*} \text{Tor}^n_{n+1}(M, C) \xrightarrow{\partial_n} \text{Tor}^n_{n-1}(M, A) \to \cdots$$

In particular, $\text{Tor}^0_R(M, N) = M \otimes_R N$.

35. Let $R$ be a commutative ring, $M$ an $R$-module, and $P_{M_*} \xrightarrow{\xi} M \to 0$ an $R$-projective resolution of $M$. Show the results in the proceeding exercise lead to $\text{Tor}^n_R(M, N) := H_n(M, P_N)$ and its long exact sequence. Show also that $\text{Tor}^n_R(M, N) = \text{Tor}^n_R(M, N)$ for all $R$-modules $M$ and $N$ and we let $\text{Tor}^n_R(M, N)$ denote both.

36. Let $I$ be a partially ordered set under $\leq$ that also satisfies for all $i, j \in I$, there exists a $k \in I$ such that $i \leq k$ and $j \leq k$; and $\{M_i\}_I$ be a collection of $R$-modules. Suppose for all $i \leq j$ in $I$, there exist $R$-homomorphisms $\theta_{ij} : M_i \to M_j$ with $\theta_{ii} = 1_{M_i}$ and

commutes whenever $i \leq j \leq k$. We call $(M_i, \theta_{ij})_I$ a directed system. Show that there exists an $R$-module $M$ and for each $i \in I$ an $R$-homomorphism $\psi_i : M_i \to M$ satisfying the following universal property: For all non-negative $i \leq j$ in $I$,
commutes, and if there exist \( R \)-homomorphisms \( \theta_i : M_i \to M \) satisfying for all \( i \leq j \) in \( I \),

\[
\begin{array}{c}
M_i \\
\theta_{ij} \\
M_j \\
\end{array} \xrightarrow{\varphi_i} \begin{array}{c}
M' \\
\varphi_j \\
\end{array}
\]

commutes, then there exists a unique \( R \)-homomorphism \( \mu : M_i \to M' \) satisfying

\[
\begin{array}{c}
M_i \\
\theta_{ij} \\
M_j \\
\end{array} \xrightarrow{\psi_i} \begin{array}{c}
M \\
\psi_j \\
\end{array} \xrightarrow{\mu} \begin{array}{c}
M' \\
\end{array}
\]

commutes for all \( i \leq j \) in \( I \). Such a \( M \) is unique up to a unique isomorphism and is called the direct limit of the \( M_i \) and denoted by \( \lim_{\rightarrow} M_i \).

37. \( R \) be a commutative ring and \((M_i, \theta_{ij})_I \) a directed system of \( R \)-modules. Show if \( A \) is an \( R \)-module, then

\[
\lim_{\rightarrow} M_i \otimes_R A \cong \lim_{\rightarrow} (M_i \otimes_R A)
\]

\[
A \otimes_R \lim_{\rightarrow} M_i \cong \lim_{\rightarrow} (A \otimes_R M_i)
\]

38. Let \( R \) be a commutative ring and \( N \) an \( R \)-module. We say that \( N \) is a flat \( R \)-module if whenever

\[
0 \to M' \xrightarrow{f} M \xrightarrow{g} M'' \to 0
\]

is exact, so is

\[
0 \to M' \otimes_R N \xrightarrow{f \otimes 1_N} M \otimes_R N \xrightarrow{g \otimes 1_N} M'' \otimes_R N \to 0.
\]

Show that if \( N \) is a projective \( R \)-module, then it is \( R \)-flat.

39. Let \( R \) be a commutative ring and \( N \) an \( R \)-module. Show that \( N \) is \( R \)-flat if and only if every finitely generated submodule of \( N \) is flat.
40. Let \( R \) be a commutative ring. If \( M \) is an \( R \)-module, let \( M^* = \text{Hom}_\mathbb{Z}(M, \mathbb{Q}/\mathbb{Z}) \). Show that

\[
0 \to A \xrightarrow{f} B \xrightarrow{g} C \to 0
\]

is an exact sequence of \( R \)-modules if and only if

\[
0 \to C \xrightarrow{g^*} B \xrightarrow{g^*} A \to 0
\]

is exact.

41. (Lambek’s Theorem) Let \( R \) be a commutative ring and \( M \) an \( R \)-module. Prove that \( M \) is \( R \)-flat if and only if \( M^* \) is \( R \)-injective.

42. Let \( R \) be a commutative ring and \( M \) an \( R \)-module. Prove that \( M \) is \( R \)-flat if and only if given any finitely generated left ideal \( \mathfrak{A} \) of \( R \) then

\[
M \otimes_R \mathfrak{A} \xrightarrow{1_M \otimes \text{inc}} M \otimes_R R \text{ is an } R\text{-monomorphism.} \]

[So identifying \( M \otimes_R R \) with \( M \), we have \( M \otimes_R \mathfrak{A} \to B \mathfrak{A} \) is an isomorphism.]

43. Let \( R \) be a commutative ring and \( M \) and \( R \)-module. Show the following are equivalent:

(i) \( M \) is \( R \)-flat.

(ii) \( \text{Tor}_n^R(M, N) = 0 \) for all \( n \geq 1 \) and all \( R \)-modules \( N \).

(iii) \( \text{Tor}_1^R(M, N) = 0 \) for all \( R \)-modules \( N \).

(iv) \( \text{Tor}_1^R(M, N) = 0 \) for all cyclic \( R \)-modules \( N \).

[Of course, the same if we reverse the variables.]
CHAPTER IX

Noetherian Rings and Modules

Although this chapter is relatively short, it includes some of the most important basic results in commutative algebra, viz. the theorems of Hilbert introduced in Section 35. It is based upon the concept of a Noetherian ring, previously introduced, and its generalization to modules. Hilbert’s theorems are necessary to begin the study of algebraic geometry. In particular, they show any set of polynomials not generating the unit ideal in \( F[t_1, \ldots, t_n] \), with \( F \) an algebraically closed field, have a common zero. A brief introduction to affine plane curves is also discussed.

39. Noetherian Modules

**Proposition 39.1.** Let \( R \) be a ring and \( M \) an \( R \)-module. Then the following are equivalent:

1. Every submodule of \( M \) is finitely generated.
2. \( M \) satisfies the ascending chain condition, i.e., if \( M_i \subset M \) are submodules and
   
   \[
   M_1 \subset M_2 \subset \cdots \subset M_n \subset \cdots,
   \]

   then there exists a positive integer \( N \) such that \( M_N = M_{N+i} \) for all \( i \geq 0 \). We say every ascending chain of submodules of \( M \) stabilizes. Equivalently, there exists no infinite chain
   
   \[
   M_1 < M_2 < \cdots < M_n < \cdots.
   \]

3. \( M \) satisfies the Maximum Principle, i.e., if \( S \) is a non-empty set of submodules of \( M \), then \( S \) contains a maximal element, that is a module \( M_0 \in S \) such that if \( M_0 \subset N \) with \( N \in S \), then \( N = M_0 \).

An \( R \)-module \( M \) satisfying any (hence all) of these equivalent conditions is called a **Noetherian** \( R \)-module or \( R \)-**Noetherian**.

**Proof.** (1) \( \Rightarrow \) (2): Let

\[
\mathcal{C} : \quad M_1 \subset M_2 \subset \cdots \subset M_n \subset \cdots
\]
be a chain of submodules of \( M \). It follows that the subset \( M' := \bigcup_{i=1}^{\infty} M_i \subset M \) is a submodule. By (1), it is finitely generated, so we can write \( M' = \sum_{i=1}^{n} Rx_i \) for some \( x_i \in M' \). By definition, \( x_i \in M_{j_i} \), some \( j_i \). Let \( s \) be the maximum of the finitely many \( j_i \)'s. Then \( M' = M_s \). It follows that \( M_s = M' = M_{s+i} \) for all \( i \geq 0 \).

(2) \( \Rightarrow \) (3): Suppose that \( S \) is a non-empty set of submodules of \( M \). Let \( M_1 \) lie in \( S \). If \( M_1 \) is not maximal, there exists an \( M_2 \in S \) with \( M_1 < M_2 \). Inductively, if \( M_i \) is not maximal, there exists an \( M_{i+1} \in S \) with \( M_i < M_{i+1} \). By the ascending chain condition, the sequence \( M_1 < M_2 < \cdots < M_i < \cdots \) must terminate.

[To prove this statement completely, note that we made a choice of \( M_2 \) to contain \( M_1 \). This really needs the Axiom of Choice. (Cf. the note in Exercise 29.22(6).)]

(3) \( \Rightarrow \) (1): Let \( N \subset M \) be a submodule and set

\[ S := \{ M_i \mid M_i \subset N \text{ is a finitely generated submodule} \}. \]

Then \( 0 \in S \) so \( S \neq \emptyset \). By assumption, there exists a maximal element \( M' \in S \). If \( N \neq M' \), then there exists \( x \in N \setminus M' \). But \( M' \) finitely generated means that \( M' + Rx \subset N \) is also finitely generated, so the submodule \( M' + Rx \) of \( N \) lies in \( S \). This contradicts the maximality of \( M' \). Hence \( N = M' \) is finitely generated. \( \square \)

**Remark 39.2.** To use Zorn’s Lemma instead of the Axiom of Choice, we proceed as follows.

(1) \( \Rightarrow \) (3): Let \( S \) is a non-empty set of submodules of \( M \) ordered by \( \subset \). If \( C \) is a chain in \( S \), then \( \bigcup_{C} N \) is a submodule of \( M \) so finitely generated by (1). It follows that there exists an \( N_0 \in C \) such that \( \bigcup_{C} N \subset N_0 \). As \( N' \subset \bigcup_{C} N \) for all \( N' \in C \), \( \bigcup_{C} N = N_0 \) lies in \( C \) and is an upper bound for \( C \). Zorn’s Lemma now gives a maximal element in \( S \). (Cf. Exercise 29.22(7).)

(2) \( \Rightarrow \) (1): Let \( M \) be generated by \( \{ x_1, x_2, \ldots \} \). We may assume that with \( Rx_1 < Rx_1 + Rx_2 < \cdots \). As this must be a finite chain by (2), the module \( M \) must be finitely generated.

(1) \( \Rightarrow \) (2) and (3) \( \Rightarrow \) (1) are as before.

**Remark 39.3.** Let \( R = F[t_1, \ldots, t_n, \ldots] \) (infinitely many \( t_i \)). Let \( M = R \) as an \( R \)-module. Then \( M \) is finitely generated since cyclic but the ideal \( (t_1, \ldots, t_n, \ldots) \) is clearly not finitely generated, so \( R \) is not a Noetherian \( R \)-module. Thus, in general, submodules of finitely generated modules need not be finitely generated.
Definition 39.4. Let $R$ be a commutative ring. We say that $R$ is a Noetherian ring if $R$ is a Noetherian $R$-module.

Note that Definition 39.4 agrees with our previous definition of a Noetherian ring.

We need another Noetherian $R$-module result. It is a crucial property about Noetherian modules.

Proposition 39.5. Let $M$ be an $R$-module and $N$ a submodule of $M$. Then $M$ is $R$-Noetherian if and only if $N$ and $M/N$ are $R$-Noetherian. In particular, if

$$0 \rightarrow M' \rightarrow M \rightarrow M'' \rightarrow 0$$

is an exact sequence of $R$-modules with two of the modules $M, M', M''$ being $R$-Noetherian, then they all are $R$-Noetherian.

Proof. $\Rightarrow$: Since $N_0 \subset N$ is a submodule, $N_0 \subset M$ is a submodule hence finitely generated – or any ascending chain in $N$ is an ascending chain in $M$. Thus $N$ is $R$-Noetherian. By the Correspondence Principle, a (countable) chain of submodules in $M/N$ has the form $M_1/N \subset M_2/N \subset \cdots$ where $N \subset M_1 \subset M_2 \subset \cdots$ is a chain of submodules of $M$. Thus there exists an $r$ such that $M_r = M_{r+j}$ for all $j \geq 0$ and hence $M_r/N = M_{r+j}/N$ for all $j \geq 0$.

$\Leftarrow$ is left as an exercise (cf. the analogous solvable groups result). □

Corollary 39.6. If $M, N$ are Noetherian $R$-modules, so is $M \amalg N$.

Proof. $(M \amalg N)/N \cong M$ and $N$ are Noetherian. □

The following result shows that the collection of Noetherian $R$-modules is a good generalization of finite dimensional vector spaces.

Theorem 39.7. Let $R$ be a Noetherian ring. If $M$ is a finitely generated $R$-module, then $M$ is $R$-Noetherian.

Proof. Suppose $M = \sum_{i=1}^{n} Rx_i$. Let $f : R^n \rightarrow M$ be the $R$-epimorphism given by $e_i \mapsto x_i$, where $\{e_1, \ldots, e_n\}$ is the standard basis for $R^n$. Since $R$ is $R$-Noetherian so is $R^n$ by Corollary 39.6, and hence so is $M \cong R^n / \ker f$ by Proposition 39.5. □

Corollary 39.8. Let $R$ be a Noetherian ring and $M$ a finitely generated $R$-module. Then there exists positive integers $m$ and $n$ and an exact sequence

$$R^n \xrightarrow{g} R^n \xrightarrow{f} M \rightarrow 0.$$
Proof. As is the proof of the theorem there exists an integer $n$ and an $R$-epimorphism $f : R^n \to M$. Since $R^n$ is finitely generated hence $R$-Noetherian, $\ker f$ is also finitely generated. Therefore, there exists an $R$-epimorphism $h : R^m \to \ker f$. Setting $g$ to be the composition of $h$ and the inclusion yields the result.

Remark 39.10. If $M$ is an $R$-module for which a sequence (39.9) exists, we say that $M$ is a finitely presented $R$-module. The corollary says that any finitely generated module over a Noetherian ring is finitely presented. This is very useful. A commutative ring in which every finitely generated module is finitely presented is called a coherent ring. Therefore, every Noetherian ring is coherent.

Proposition 39.11. Suppose that $f : R \to S$ is a ring epimorphism of commutative rings. If $R$ is a Noetherian ring, then $S$ is also a Noetherian ring.

Proof. Let $\mathfrak{a} \subset S$ be an ideal. Then $f^{-1}(\mathfrak{a}) \subset R$ is an ideal hence finitely generated. Thus $\mathfrak{a} = f(f^{-1}(\mathfrak{a}))$ is finitely generated. □

Exercises 39.12.
1. Complete the proof of Proposition 39.5.
2. Let $M$ be a Noetherian $R$-module. Show that any surjective $R$-endomorphism $f : M \to M$ is an isomorphism.
3. Let $M$ be an $R$-module. Prove the following are equivalent:
   (i) $M$ satisfies DCC (the descending chain condition), i.e., if $M_i \subset M$ are submodules and
   \[ M_1 \supset M_2 \supset \cdots \supset M_n \supset \cdots \]
   then there exists a positive integer $N$ such that $M_N = M_{N+i}$ for all $i \geq 0$.
   (ii) $M$ satisfies the Minimum Principle, i.e., if $S$ is a non-empty set of submodules of $M$, then $S$ contains a minimal element, that is a module $M_0 \in S$ such that if $M_0 \supset N$ with $N \in S$, then $N = M_0$.
If $M$ satisfies one of these two equivalent conditions, we say that $M$ is an Artinian $R$-module.
[One needs the Axiom of Choice to show DCC implies the Minimum Principle.]
Also show that if $R$ is a field then an $R$-module $V$ is Artinian if and only if $V$ is finite dimensional.
4. Let
\[ 0 \to M' \to M \to M'' \to 0 \]
be an exact sequence of \( R \)-modules. Show that if two of the modules \( M, M', M'' \) are \( R \)-Artinian, then they all are \( R \)-Artinian.

5. Let \( R \) be a Noetherian ring, \( \mathfrak{A} \) an ideal in \( R \), \( A \) a finitely generated \( R \)-module, and \( B \) a submodule of \( A \). Suppose that \( C \) is a submodule of \( A \) that contains \( \mathfrak{A}B \) and is maximal with respect to the property that \( C \cap B = \mathfrak{A}B \). Let \( x \) be an element of \( \mathfrak{A} \). Show all of the following:

\( i \) The chain of submodules of \( A \), \( D_m = \{ a \in A \mid x^m a \in C \text{ for all } a \in A \} \), \( m \in \mathbb{Z}^+ \), stabilizes.

\( ii \) There exists an integer \( n \) such that \( (x^n A + C) \cap B = \mathfrak{A}B \).

\( iii \) \( \mathfrak{A}^n A \subset C \) for some \( n \).

\( iv \) (Krull Intersection Theorem) If \( B = \bigcap_{i=0}^{\infty} \mathfrak{A}^i A \), then \( \mathfrak{A}B = B \).

6. Let \( R \) be a commutative ring, \( \mathfrak{A} \) an ideal in \( R \), \( M \) an \( R \)-module generated by \( n \) elements, and \( x \) an element of \( R \) satisfying \( xM \subset M \). Show that \( (x^n + y)M = 0 \) for some \( y \in \mathfrak{A} \). In particular, if \( \mathfrak{A}M = M \), then \( (1 + y)M = 0 \) for some \( y \in \mathfrak{A} \).

7. Let \( R \) be a Noetherian ring. Using the previous two exercises, show the following:

\( i \) Suppose that \( R \) is a domain and \( \mathfrak{A} \subset R \) be an ideal. Let \( M \) be a finitely generated \( R \)-module satisfying \( \text{ann}_R(m) = 0 \) for all \( m \in M \). (We say that \( M \) is \( R \)-torsion-free.) Then \( \bigcap_{i=0}^{\infty} \mathfrak{A}^i M = 0 \).

\( ii \) Let \( R \) be a Noetherian ring \( \mathfrak{A} = \bigcap_{\text{Max}(R)} \mathfrak{m} \), the Jacobson radical of \( R \), (cf. Exercise 27.19(8)), and \( M \) a finitely generated \( R \)-module. Then \( \bigcap_{i=0}^{\infty} \mathfrak{A}^i M = 0 \).

40. Hilbert’s Theorems

In this section, we prove the theorems that we discussed in Section 35. These theorems form the foundation of affine algebraic geometry.

**Theorem 40.1.** (Hilbert Basis Theorem) If \( R \) is a Noetherian ring so is the ring \( R[t_1, \ldots, t_n] \).

**Proof.** By induction on \( n \), it suffices to show that \( R[t] \) is Noetherian. Let \( 0 \neq \mathfrak{B} \subset R[t] \) be an ideal. We must show that \( \mathfrak{B} \) is finitely generated. Let

\[ \mathfrak{A} = \{ r \in R | r = \text{lead } f, f \in \mathfrak{B} \} \]

(Recall lead \( f \) is the leading coefficient of \( f \).)

We first show that \( \mathfrak{A} \) is an ideal. (Note that \( 0 \notin \mathfrak{A} \) as \( 0 \notin \mathfrak{B} \).) Let \( a, b \in \mathfrak{A} \) and \( r \in R \). To show \( ra + b \) lies in \( \mathfrak{A} \). We may assume that \( a \) and \( b \) are nonzero. Choose \( f, g \in \mathfrak{B} \) say of degrees \( m \) and \( n \)
respectively satisfying \( \text{lead } f = a \) and \( \text{lead } g = b \). Set \( h = rt^n f + tm^g \) in \( \mathfrak{B} \). Then \( ra + b = \text{lead } h \) proving \( \mathfrak{A} \) is an ideal. As \( R \) is Noetherian, \( \mathfrak{A} = (a_1, \ldots, a_n) \) for some nonzero \( a_1, \ldots, a_n \in \mathfrak{A} \) with \( n \in \mathbb{Z}^+ \). Choose \( f_{d_i} \) in \( \mathfrak{B} \) such that \( a_i = \text{lead } f_{d_i} \) and \( \deg f_{d_i} = d_i \) for \( i = 1, \ldots, n \). Let \( \mathfrak{B}_0 = (f_{d_1}, \ldots, f_{d_n}) \), an ideal in \( R[t] \), and \( N = \max\{d_1, \ldots, d_n\} \).

Let \( f \in \mathfrak{B} \) with \( \text{lead } f = a \) and \( \deg f = d \). Suppose that \( d > N \). There exist \( r_i \in R \) satisfying \( a = \sum_{i=1}^n r_i a_i \), hence \( f - \sum_{i=1}^n r_i t^{d_i - d} f_{d_i} \) lies in \( \mathfrak{B} \) and has degree less than \( d \). It follows by induction that there exists a \( g \in \mathfrak{B}_0 \) such that \( f - g \) lies in \( \mathfrak{B} \) with \( \deg(f - g) \leq N \). As the \( R \)-module \( M := \sum_{i=0}^N R[t]^i \) is finitely generated and \( R \) is Noetherian, \( M \) is a Noetherian \( R \)-module. In particular, the submodule \( M_0 = \{f \in \mathfrak{B} \mid \deg f \leq N\} \) is finitely generated. If \( M_0 = \sum_{i=0}^m R[t]g_i \), then \( \mathfrak{B} = \mathfrak{B}_0 + \sum_{i=0}^m R[t]g_i \) is finitely generated. \( \square \)

It is worth noting the similarity of this proof and that for the general division algorithm. Of course, the division algorithm is a much stronger property, but the idea of a proof often generalizes.

**Remarks 40.2.** 1. Noetherian domains need not be UFD’s, as \( \mathbb{Z}[\sqrt{-5}] \) is an example of such a domain.

2. If \( F \) is a field, \( F[t_1, \ldots, t_n, \ldots] \) is a UFD but is not Noetherian.

**Definition 40.3.** Let \( R \subset S \) be commutative rings. We say that \( S \) is a *finitely generated commutative \( R \)-algebra* (or an *affine \( R \)-algebra* when \( R \) is a field) if there exist \( x_1, \ldots, x_n \) in \( S \) satisfying \( S = R[x_1, \ldots, x_n] \) as rings. (Recall that \( R[x_1, \ldots, x_n] \) is the image of the evaluation map \( e_{x_1, \ldots, x_n} : R[t_1, \ldots, t_n] \to S \).)

**Remarks 40.4.** Let \( R \) be a commutative ring.

1. \( R[t] \) is a finitely generated commutative \( R \)-algebra but definitely not a finitely generated \( R \)-module. (Why?)

2. Let \( S \) be a commutative ring with \( R \) a subring. If \( S \) is a finitely generated \( R \)-module, then \( S \) is a finitely generated commutative \( R \)-algebra.

3. If \( S \) is a finitely generated commutative \( R \)-algebra, it is not in general a finitely generated \( R \)-module, a much stronger condition. (Cf. (1).)

4. If \( T \) is a finitely generated commutative \( S \)-algebra and \( S \) is a finitely generated commutative \( R \)-algebra, then \( T \) is a finitely generated commutative \( R \)-algebra.

5. If \( T \) is a commutative ring with \( S \) a subring such that \( T \) is a finitely generated \( S \)-module and \( S \) is a finitely generated commutative \( R \)-algebra, then \( T \) is a finitely generated commutative \( R \)-algebra.
Proposition 40.5. Let \( R \) be a commutative ring and \( S \) a finitely generated commutative \( R \)-algebra. If \( R \) is Noetherian so is \( S \).

**Proof.** Let \( S = R[x_1, \ldots, x_n] \). Since \( R[t_1, \ldots, t_n] \) is Noetherian and we have a ring epimorphism \( R[t_1, \ldots, t_n] \to R[x_1, \ldots, x_n] \) via \( f(t_1, \ldots, t_n) \mapsto f(x_1, \ldots, x_n) \), all ideals of \( S \) are finitely generated by the Correspondence Principle. □

Remark 40.6. Let \( R \) be a commutative ring. The proof (or definition) shows that \( S \) is a finitely generated commutative \( R \)-algebra if and only if there exists an ring epimorphism \( R[t_1, \ldots, t_n] \to S \) fixing \( R \). Our definition of a finitely generated commutative \( R \)-algebra is not the usual one. In general, one does not assume that \( R \subset S \), hence only that there exists a surjective ring homomorphism \( R[t_1, \ldots, t_n] \to S \) for some \( n \). If \( R \) is a field, we can always assume that \( R \subset S \), as the above map must be monic when restricted to \( R \) if \( R \) is a field. This is the case of interest here, except for the following lemma.

Lemma 40.7. (Artin-Tate) Let \( T \) be a commutative ring and \( R \subset S \) be subrings of \( T \). Suppose that \( R \) is Noetherian and \( T \) is a finitely generated commutative \( R \)-algebra. Suppose that as an \( S \)-module \( T \) is finitely generated. Then \( S \) is a finitely generated commutative \( R \)-algebra. In particular, \( S \) is a Noetherian ring.

**Proof.** Let \( T = R[x_1, \ldots, x_n] \) as a commutative \( R \)-algebra for some \( x_i \in T \), some \( n \) and \( T = \sum_{i=1}^{m} S y_i \) as an \( S \)-module for some \( y_i \in T \), some \( m \). Then for all \( i = 1, \ldots, n \) and \( p, q = 1, \ldots, m \), we have equations:

1. \( x_i = \sum_{j=1}^{m} a_{ij} y_j \) for some \( a_{ij} \in S \) and
2. \( y_p y_q = \sum_{k=1}^{m} b_{pqk} y_k \) for some \( b_{pqk} \in S \) (since \( T \) is a ring).

Let \( S_0 = R[a_{ij}, b_{pqk} \mid i = 1, \ldots, n \text{ and } j, p, q, k = 1, \ldots, m] \), a finitely generated \( R \)-algebra. Then \( R \subset S_0 \subset S \subset T \).

**Claim.** \( T \) is a finitely generated \( S_0 \)-module.

Let \( f \in T \), so \( f = \sum c_{i_1, \ldots, i_n} x_1^{i_1} \cdots x_n^{i_n} \) for some \( c_{i_1, \ldots, i_n} \in R \). Applying properties (i) and (ii) repeatedly shows that \( f \in \sum_{i=1}^{m} S_0 y_i \). Consequently, \( T = S_0 y_1 + \cdots + S_0 y_m \) as claimed.
Since $S_0$ is a finitely generated commutative $R$-algebra, it is Noetherian by Corollary 40.5. Thus $T$, being a finitely generated $S_0$-module, is a Noetherian $S_0$-module. Consequently, $S \subset T$ is a finitely generated $S_0$-module. It follows immediately that $S$ is a finitely generated commutative $R$-algebra. □

Let $R$ be an affine $F$-algebra, say $R = F[x_1, \ldots , x_n]$. If $R$ is a domain, we write the quotient field of $R$ by $F(x_1, \ldots , x_n)$. We shall need the following two lemmas about fields that we leave as easy exercises.

**Lemma 40.8.** Suppose that $L \subset K \subset M$ are fields. Then $M$ is a finite dimensional vector space over $L$ if and only if both $M$ is a finite dimensional vector space over $K$ and $K$ is a finite dimensional vector space over $L$.

**Lemma 40.9.** Let $L \subset M$ be fields. Suppose that $x \in M$ has the property that $L[x]$ is not a finite dimensional vector space over $L$. Then $L[x] \cong L[t]$ as rings and $L(x) \cong L(t)$ as fields.

The clever Artin-Tate Lemma allows us to give an elementary proof of the following result that is essentially a form of the Hilbert Nullstellensatz.

**Lemma 40.10.** (Zariski’s Lemma) Let $F$ be a field and $E$ a field containing $F$ such that $E$ is an affine $F$-algebra. Then $E$ is a finite dimensional vector space over $F$ (i.e., a finitely generated $F$-module).

**Proof.** Suppose that $E = F[x_1, \ldots , x_m]$. Since $E$ is a field, $E = F(x_1, \ldots , x_m)$. Suppose that $E$ is not a finite dimensional vector space over $F$. Using Lemma 40.8, we see that $F(x_i)$ is an infinite dimensional vector space over $F$ for some $i$. Relabeling, we may assume $i = 1$. Continuing in this way, we see that after relabeling the $x_i$, we may assume that $F(x_1, \ldots , x_i)$ is not a finite dimensional vector space over $F(x_1, \ldots x_{i-1})$ for $1 \leq i \leq r$, some $r$, and $E$ is a finite dimensional vector space over $F(x_1, \ldots , x_r)$.

Let $K = F(x_1, \ldots , x_r) \subset E$. We have $E = K(x_{r+1}, \ldots , x_m)$ is a finitely generated $K$-module and $F$ is a Noetherian ring (since a field). Thus the Artin-Tate Lemma implies that $K$ is a finitely generated $F$-algebra. Write $K = F[y_1, \ldots , y_n]$ for some $y_i \in K$. By Lemma 40.9, it follows that and $F[x_1, \ldots , x_r] \cong F[t_1, \ldots , t_r]$ and $K \cong F(t_1, \ldots , t_r)$ as rings. Thus we can write

$$y_i = \frac{f_i(x_1, \ldots , x_r)}{g_i(x_1, \ldots , x_r)}$$

for some $f_i, g_i \in F[t_1, \ldots , t_r]$, $g_i \neq 0$, $1 \leq i \leq n$. Let $g = g_1 \cdots g_n$ in $F[t_1, \ldots , t_r]$. Thus $g(x_1, \ldots , x_r) \in F[x_1, \ldots , x_r]$. We know the UFD
\[ F[t_1, \ldots, t_r] \text{ contains infinitely many non-associative irreducibles hence prime elements by Exercise 33.19(9). In particular, there exists an irreducible } f \in F[t_1, \ldots, t_r] \text{ such that } f \not| g. \] Thus \( f(x_1, \ldots, x_r) \not| g(x_1, \ldots, x_r) \) in \( F[x_1, \ldots, x_r] \). As \( K \) is a field,

\[
\frac{1}{f(x_1, \ldots, x_r)} \in K = F[y_1, \ldots, y_n] = F\left[ f_1(x_1, \ldots, x_r), \ldots, f_n(x_1, \ldots, x_r), \right].
\]

This leads to an equation for \( 1/f(x_1, \ldots, x_r) \). Choosing an appropriate \( N \geq 0 \) to clear denominators, we see that we have an equation

\[
\frac{g(x_1, \ldots, x_r)^N}{f(x_1, \ldots, x_r)} \in F[f_1(x_1, \ldots, x_r), \ldots, f_n(x_1, \ldots, x_r)].
\]

Then \( \frac{g(x_1, \ldots, x_r)^N}{f(x_1, \ldots, x_r)} \) lies in \( F[x_1, \ldots, x_r] \), i.e., \( f|g^N \) in \( F[t_1, \ldots, t_r] \), a contradiction.

Recall from §35, if \( R = F[t_1, \ldots, t_n] \) and \( \mathfrak{A} \subset R \) is an ideal, the \textit{affine variety} of \( \mathfrak{A} \) in \( F^n \) is defined by

\[
Z_F(\mathfrak{A}) = \{ \mathfrak{a} = (a_1, \ldots, a_n) \in F^n \mid f(\mathfrak{a}) = 0 \text{ for all } f \in \mathfrak{A} \}.
\]

By the Hilbert Basis Theorem, there exist \( f_1, \ldots, f_r \) in \( F[t_1, \ldots, t_n] \) satisfying \( \mathfrak{A} = (f_1, \ldots, f_r) \). We then also write \( Z_F(f_1, \ldots, f_r) \) for \( Z_F(\mathfrak{A}) \).

In particular, \( \mathfrak{a} \) lies in \( Z_F(\mathfrak{A}) \) if and only if \( f_i(\mathfrak{a}) = 0 \) for \( i = 1, \ldots, r \), as any element in \( \mathfrak{A} \) is an \( R \)-linear combination of the \( f_i \)'s.

**Theorem 40.11.** (Hilbert Nullstellensatz) (Weak Form) \textit{Suppose that } \( F \) \textit{ is an algebraically closed field, } \( R = F[t_1, \ldots, t_n] \), \textit{ and } \( \mathfrak{A} = (f_1, \ldots, f_r) \) \textit{ is an ideal in } \( R \). \textit{Then } \( Z_F(\mathfrak{A}) \) \textit{ is the empty set if and only if } \( \mathfrak{A} \) \textit{ is the unit ideal. Moreover, if } \( \mathfrak{A} \subset R \), \textit{ then there exists a point } \( \mathfrak{a} \in F^n \) \textit{ satisfying } \( f_1(\mathfrak{a}) = 0, \ldots, f_r(\mathfrak{a}) = 0 \).

**Proof.** Certainly if \( \mathfrak{A} = R \), then there exists no \( \mathfrak{a} \in F^n \) such that the element 1 in \( R \) evaluated at \( \mathfrak{a} \) takes the value zero, so we need only show if \( \mathfrak{A} < R \), then \( Z_F(\mathfrak{A}) \) is not empty.

So suppose that \( \mathfrak{A} < R \). Then there exists a maximal ideal \( \mathfrak{m} < R \) satisfying \( \mathfrak{A} \subset \mathfrak{m} \). Let \( \frac{1}{\mathfrak{m}} : R \to R/\mathfrak{m} \) be the canonical ring epimorphism.

Set \( E = R/\mathfrak{m} = \overline{\mathfrak{m}}[t_1, \ldots, t_n] \). As \( R \) is an affine \( F \)-algebra, so is \( E \). By Zariski’s Lemma, \( E \) is a finite dimensional \( F \)-vector space.

**Claim.** \( E = F \):

Indeed let \( x \in E \). Then \( F[x] \) is a finite dimensional \( F \)-vector space, so \( \{1, x, x^2, \ldots, x^N\} \) must be linear dependent over \( F \) for some \( N \), i.e., \( x \) is a root of some non-zero polynomial \( f \in F[t] \). But any such \( f \) factors completely over \( F \), as \( F \) is algebraically closed. Therefore, \( E = F \).
Consequently the point \( t = (t_1, \ldots, t_n) \) in \( E^n = F^n \) lies in \( Z_F(\mathfrak{A}) \), since \( \mathfrak{A} \subset \mathfrak{m} \), i.e., \( \mathfrak{A} = \mathfrak{m} = 0 \) □

Recall if \( F \) is a field and \( \mathfrak{m}_x \) is the ideal \((t_1 - x_1, \ldots, t_n - x_n)\) with \( x = (x_1, \ldots, x_n) \) in \( F^n \), then \( \mathfrak{m}_x \) is a maximal ideal and Equation 35.2 says

\[
\mathfrak{A} \subset \mathfrak{m}_x \text{ if and only if } x \in Z_F(\mathfrak{A}).
\]

for an ideal \( \mathfrak{A} \) in \( F[t_1, \ldots, t_n] \).

The Weak Form of the Hilbert Nullstellensatz says that the analogue of Exercise 25.21(13) that the maximal ideals in the ring of continuous real-valued functions on a finite closed (or compact) set are in one to one correspondence with the points of the set is valid over an algebraically closed field.

**Corollary 40.12.** Let \( F \) be an algebraically closed field and \( \mathfrak{m} \) a maximal ideal in \( F[t_1, \ldots, t_n] \). Then there exists an element \( x \) in \( F^n \) such that \( \mathfrak{m} = \mathfrak{m}_x \).

**Proof.** By the Weak Nullstellensatz, there exists a point \( x \) in \( Z_F(\mathfrak{m}) \). Let \( - : F[t_1, \ldots, t_n] \to F[t_1, \ldots, t_n]/\mathfrak{m}_x \) be the canonical epimorphism. It takes \( t_i \to \overline{t}_i = x_i \) for each \( i \). By (*) we have \( \mathfrak{m} \subset \mathfrak{m}_x \). So \( \mathfrak{m} = \mathfrak{m}_x \) as \( \mathfrak{m} \) is a maximal ideal.

Note that the corollary says if \( F \) is algebraically closed, then there is a bijection between points in \( F^n \) and maximal ideals in \( F[t_1, \ldots, t_n] \) and the points in \( Z_F(\mathfrak{A}) \) bijective with the maximal ideals containing \( \mathfrak{A} \) in \( F[t_1, \ldots, t_n] \). (Cf. Exercise 25.21(13).)

The geometric interpretation of the Weak Hilbert Nullstellensatz is that the intersection of the “hyperplanes”

\[ f_1 = 0, \ldots, f_r = 0 \quad \text{in} \quad F^n. \]

always contains a common point over an algebraically closed field \( F \), unless the \( f_i \) generate the unit ideal in \( F[t_1, \ldots, t_n] \).

Recall also that \( F \) being algebraically closed is essential. Indeed \( f(t_1, \ldots, t_n) = t_1^2 + \cdots + t_n^2 + 1 \) has no solution in \( \mathbb{R}^n \) yet \( (f) < \mathbb{R}[t_1, \ldots, t_n] \). Can you state what the above argument shows when \( F \) is not algebraically closed?

The Weak Hilbert Nullstellensatz says that \( f_1, \ldots, f_r \) determine \( Z_F(f_1, \ldots, f_r) \) when \( F \) is algebraically closed. The natural question is what polynomials in \( F[t_1, \ldots, t_n] \) does \( Z_F(f_1, \ldots, f_r) \) determine when \( F \) is algebraically closed, i.e., what polynomials in \( F[t_1, \ldots, t_n] \) have zeros at every point in \( Z_F(f_1, \ldots, f_r) \). Certainly elements in \( \mathfrak{A} = (f_1, \ldots, f_r) \) have this property. However, a point in \( F^n \) is a zero of a polynomial \( f \)
in the domain $F[t_1, \ldots, t_n]$ if and only if it is a zero of $f^m$ for any $m > 0$ if and only if it is a zero of $f^m$ for some $m > 0$. This means that $Z_F(\mathfrak{A}) = Z_F(\sqrt{\mathfrak{A}})$ for any ideal $\mathfrak{A}$, i.e., at a minimum, every element in the radical ideal $\sqrt{\mathfrak{A}} := \{x \mid x^n \in F[t_1, \ldots, t_n] \text{ for some positive integer } n\}$ of $\mathfrak{A}$ has every point in $Z_F(\mathfrak{A})$ a zero. The Strong Hilbert Nullstellensatz says, when $F$ is algebraically closed, this is the totality of such polynomials.

**Theorem 40.13.** (Hilbert Nullstellensatz) (Strong Form) Suppose that $F$ be an algebraically closed field and $R = F[t_1, \ldots, t_n]$. Let $f, f_1, \ldots, f_r$ be elements in $R$ and $\mathfrak{A} = (f_1, \ldots, f_r) \subset R$. Suppose that $f(a) = 0$ for all $a \in Z_F(\mathfrak{A})$. Then there exists an integer $m$ such that $f^m \in \mathfrak{A}$, i.e., $f \in \sqrt{\mathfrak{A}}$. In particular, if $\mathfrak{A}$ is a prime ideal, then $f \in \mathfrak{A}$.

**Proof.** (Rabinowitch Trick). We may assume that $f$ is nonzero. Let $S = R[t]$. Define the ideal $\mathfrak{B}$ in $S$ by $\mathfrak{B} = (f_1, \ldots, f_r, 1 - tf) \subset S$. If $\mathfrak{B} < S$, then there exists a point $(a_1, \ldots, a_{n+1}) \in Z_F(\mathfrak{B})$. Thus $f_i(a_1, \ldots, a_n) = 0$ for all $i$ and $1 - a_{n+1}f(a_1, \ldots, a_n) = 0$. In particular, $(a_1, \ldots, a_n)$ lies in $Z_F(\mathfrak{A})$. By hypothesis, this means that $f(a_1, \ldots, a_n) = 0$ which in turn implies that $1 = 0$, a contradiction. Thus $\mathfrak{B} = S$. So we can write

$$1 = \sum_{i=1}^r g_i f_i + g \cdot (1 - tf)$$

for some $g, g_i \in S$. Substituting $1/f$ for $t$ and clearing denominators yields the result. \hfill $\square$

If $R$ is a commutative ring, and $\mathfrak{A}$ an ideal in $R$, then applying the Correspondence Principle to the canonical epimorphism $- : R \to R/\mathfrak{A}$ shows that $\sqrt{\mathfrak{A}}$ is the intersection of all prime ideals containing $\mathfrak{A}$. The strong form of the Nullstellensatz has an algebraic version that says for any field $F$, if $R$ is an affine $F$-algebra and $\mathfrak{A}$ an ideal in $R$, then

$$\sqrt{\mathfrak{A}} = \bigcap_{m \text{ a maximal ideal}} m.$$ 

We do not prove this here [The proof can be found below in Theorem 94.11 using localization techniques.] Instead we indicate what is going on (also without proof) assuming the reader knows what a topology is.

Let $R$ be a nonzero commutative ring and set

$$\text{Spec}(R) := \{p \mid p < R \text{ a prime ideal}\},$$

called the *Spectrum* of $R$. If $T$ is a subset of $R$, define

$$V_R(T) := \{p \mid p \in \text{Spec}(R) \text{ with } T \subset p\}$$
called a variety. Note that $V_R(T) = V_R(\langle T \rangle)$. We leave the following as an exercise.

**Lemma 40.14.** Let $R$ be a commutative ring, $\mathfrak{A}$, $\mathfrak{B}$, and $\mathfrak{A}_i, i \in I$, be ideals in $R$. Then

1. If $\mathfrak{A} \subset \mathfrak{B}$, then $V_R(\mathfrak{B}) \subset V_R(\mathfrak{A})$.
2. $V_R(\emptyset) = \Spec(R)$
3. $V_R(R) = \emptyset$.
4. $V_R(\sum \mathfrak{A}_i) = \bigcap_i V_R(\mathfrak{A}_i)$.
5. $V_R(\mathfrak{A}\mathfrak{B}) = V_R(\mathfrak{A} \cap \mathfrak{B}) = V_R(\mathfrak{A}) \cup V_R(\mathfrak{B})$.
6. $V_R(\mathfrak{A}) = V_R(\sqrt{\mathfrak{A}})$.

Lemma (2) – (4) means that the collection

$$ C := \{V_R(T) \mid T \subset R\} $$

forms a system of closed sets for a topology on $\Spec(R)$ called the Zariski topology.

One shows that the set

$$ \Max(R) := \{m \mid m < R \text{ a maximal ideal}\} $$

of maximal ideals in $R$ is precisely the set of closed points in $\Spec(R)$, i.e., a closed set with precisely one element. If $F$ is a field, then the collection

$$ Z := \{Z_F(T) \mid T \subset F[t_1, \ldots, t_n]\} $$

also forms a system of closed sets for a topology of $F^n$ by Exercise 35.5 called the Zariski topology for $F^n$. If $F$ is algebraically closed, the weak form of the Nullstellensatz says if $\Max(R)$ is given the induced topology, then

$$ \Max(F[t_1, \ldots, t_n]) \to F^n \text{ given by } m \mapsto z $$

is a homeomorphism and the strong form of the Nullstellensatz says that $\Max(F[t_1, \ldots, t_n])$ is dense in $\Spec(F[t_1, \ldots, t_n])$, i.e., every non-empty open set in $\Spec(R)$ intersects $\Max(R)$ nontrivially. This means that the closed points determine varieties. This last statement generalizes to $\Max(R)$ is dense in $\Spec(R)$ for any affine $F$-algebra $R$. (Cf. with (*) above.)

**Exercises 40.15.**
1. Prove that if $R$ is a Noetherian ring so is the formal power series $R[[t]]$.
2. Prove Lemma 40.8. If $M$ is a finite dimensional vector space over $L$, then $\dim_L M = \dim_k M \dim_L K$? [Compare this to the generalized version of Lagrange’s Theorem for groups.]
3. Prove Lemma 40.9.

4. Let $F$ be a field and $\varphi: A \to B$ be a ring homomorphism of $F$-algebras such that $\varphi|_F = 1_A$. Prove if $m$ is a maximal ideal in $B$, then $\varphi^{-1}(m)$ is a maximal ideal in $A$.

5. Prove if $R$ is a commutative ring, then $\sqrt{\mathfrak{A}}$ is the intersection of all prime ideals containing $\mathfrak{A}$.


7. Let $R$ be a commutative ring. Prove that Spec$(R)$ contains minimal elements under $\subset$.

41. Addendum: Affine Plane Curves

Let $F$ be a field and $\mathfrak{A}$ an ideal in $F[t_1, \ldots, t_n]$. Then the affine variety $Z_F(\mathfrak{A}) = Z_F(\sqrt{\mathfrak{A}})$ (cf. Exercise 35.5). Note if $p$ is a prime ideal then $p = \sqrt{p}$. One can show:

**Fact 41.1.** Let $F$ be a field and $\mathfrak{A} < F[t_1, \ldots, t_n]$ an ideal. Then

(*) \[ Z_F(\mathfrak{A}) = Z_F(p_1) \cup \cdots \cup Z_F(p_n) \]

where $p_1, \ldots, p_n$ are the finitely many prime ideals in $F[t_1, \ldots, t_n]$ that minimally contain $\mathfrak{A}$ (i.e., if $\mathfrak{P}$ is a prime ideal of $F[t_1, \ldots, t_n]$ satisfying $\mathfrak{A} \subset \mathfrak{P}$, then there exists an $i$ such that $\mathfrak{A} \subset p_i \subset \mathfrak{P}$).

We know that $F[t_1, \ldots, t_n]$ is a Noetherian ring. To show that there are finitely many such $p_i$ in $F[t_1, \ldots, t_n]$ in Fact 41.1 uses this (cf. Exercise 29.22(15)). This reduces the study of $Z_F(\mathfrak{A})$ to $Z_F(p)$ with $p$ a prime ideal in $F[t_1, \ldots, t_n]$. Such a variety is called irreducible. The decomposition (*) for $Z_F(\mathfrak{A})$ is called the irreducible decomposition of $Z_F(\mathfrak{A})$ (cf. Exercise 29.22(15)) and the $Z_F(p_i)$ are called the irreducible components of $Z_F(\mathfrak{A})$. As $F[t_1, \ldots, t_n]$, is a UFD, it also turns out that every irreducible affine variety of “codimension one” in $F^n$ is $Z_F(f)$, for some prime element $f$. [This essentially follows from the fact that every nonzero prime ideal in a UFD contains a prime element together with Exercise 30.3(3)] and “dimension theory”, that we have not and will not develop that shows how this defines codimension algebraically.]

We shall need to use Exercise 34.12(3ii) in the development below. As we are interested in affine varieties in $F^2$, it is convenient to use $X, Y$ as our variables instead of $t_1, t_2$.

**Example 41.2.** Let $F$ be a field and $\mathfrak{A} = (XY)$ in $F[X, Y]$. Then $Z_F(\mathfrak{A}) = Z_F(X) \cup Z_F(Y)$ is an irreducible decomposition. Here $Z_F(X)$ is the $Y$-axis defined by $X = 0$ and $Z_F(Y)$ is the $X$-axis defined by $Y = 0$. The origin is the only point in $Z_F(X, Y)$. As $Z_F(X, Y) \subset Z_F(XY)$
(cf. Exercise 35.5), we see that irreducible components can intersect nontrivially. If \( f \in F[X,Y] \) is a non-constant polynomial then \( Z_F(f) \) is called an affine (plane) curve and irreducible if \( f \) is an irreducible polynomial as \( Z_F(f) \) is the locus of \( f = 0 \) in \( F^2 \).

**Proposition 41.3.** Let \( R \) be a UFD and \( f \) and \( g \) non-constant polynomials in \( R[X,Y] \) having no non-constant common factor in \( R[X,Y] \). Then

\[
Z_R(f) \cap Z_R(g) := \{(a, b) \in R^2 \mid f(a, b) = 0 = g(a, b)\}
\]

is a finite set in \( R^2 \).

**Proof.** Let \( F \) be the quotient field of \( R \) and \( K \) the quotient field of \( F[Y] \). We view \( f \) and \( g \) in \( K \), and Exercise 34.12(3ii), \( f \) and \( g \) have no common factor in \( K[Y] \), i.e., they are relatively prime in the PID \( K[Y] \). Clearing the denominators of the \( K \)-coefficients leads to an equation \( 0 \neq h = w f + z g \) in \( R[X,Y] \), with \( w \) and \( z \) in \( R[X,Y] \) and \( h \) in \( R[X] \). Any common solution \( (a, b) \) in \( Z_F(g) \cap Z_F(g) \) then satisfies \( h(a) = h(a, b) = 0 \) in \( R^2 \). Since \( h \) has finitely many roots in the domain \( R \), there exist finitely many \( a \) satisfying \( (a, b) \in Z_F(g) \cap Z_F(g) \). Applying the same argument with the variables reversed shows that \( Z_F(g) \cap Z_F(g) \) is a finite set. \( \Box \)

**Remark 41.4.** (Cf. Exercise 35.5). Let \( F \) be a field, \( \mathfrak{A} \) and \( \mathfrak{B} \) ideals in \( F[X,Y] \). Then

\[
Z_F(\mathfrak{A} \cap \mathfrak{B}) = Z_F(\mathfrak{A}\mathfrak{B}) = Z_F(\mathfrak{A}) \cup Z_F(\mathfrak{B}).
\]

The remark, reduces our study to ideals \( \mathfrak{A} = (f_1, \ldots, f_r) \) in \( F[X,Y] \) with a gcd of the \( f_1, \ldots, f_r \) equal to one, i.e., with no common factors.

If \( h \) is a polynomial in \( R[X,Y] \), \( R \) a ring, define the total degree \( \deg h \) of \( h \) to be

\[
\deg h := \max\{i + j \mid aX^iY^j \text{ a monomial in } h \text{ with } a \text{ nonzero}\}.
\]

The proposition can be improved to

**Fact 41.5.** (Bezout’s Lemma) Let \( F \) be a field and \( f \) and \( g \) two non-constant polynomials in \( F[X,Y] \) having no non-constant common factor in \( F[X,Y] \). Then \( |Z_F(f) \cap Z_F(g)| \leq \deg f \deg g \).

Our argument does not show this or give any bounds.

The proposition, in some sense, reduces our study of affine plane curves over a field to \( Z_F(f) \) with \( f \) an irreducible polynomial in \( F[X,Y] \) as \( Z_F(f) = Z_F(f^r) \) for any positive integer \( r \).
If \( f = f(X,Y) \) is a polynomial in two variables, we can view \( f \) as a polynomial in \( Y \) over \( F[X] \) or its quotient field \( F(X) \), say the degree in \( Y \) of \( f \) is \( n \). If we evaluate \( X \) at \( x \) in \( F \), then \( f(x,Y) \) is a polynomial of degree at most \( n \). We want to show if \( F \) is algebraically closed of characteristic zero, that for almost all such \( x \), the polynomial in the variable \( Y \), \( f(x,Y) \) is precisely of degree \( n \) and moreover has no multiple roots in \( F \), where \( a \) is a multiple root of \( f(x,Y) \) in \( F \) if \( (Y - a)^2 \mid f(x,Y) \) in \( F[Y] \) for any root \( a \) of \( f(x,Y) \) in \( F \).

**Definition 41.6.** Let \( f \) be a polynomial in \( F[X,Y] \) of degree \( \text{deg}_Y f \) of \( f \) in the variable \( Y \) satisfying \( \text{deg}_Y f = n > 0 \). An element \( x \) in \( F \) is called a regular value of \( f \) if \( f(x,Y) \) has \( n \) distinct roots in \( F \).

**Remark 41.7.** Let \( F \) be a field of characteristic zero and \( x \) a regular value of \( f \) in \( F[X,Y] \) of degree \( n \) in \( Y \). We can write \( f = \sum_{i=0}^{n} q_i(X)Y^i \) with \( q_i(X) \) polynomials in \( F[X] \) and \( q_n(X) \) nonzero. As \( x \) is a regular value, \( q_n(x) \) is nonzero and \( f(x,Y) \) has \( n \) distinct roots. In particular none of the roots of \( f(x,Y) \) are roots of the formal derivative \( \frac{\partial f}{\partial Y}(x,Y) \) by Exercise 33.19(3ii) or by Exercise 41.10(4).

**Proposition 41.8.** Let \( F \) be an algebraically closed field of characteristic zero, e.g., \( F = \mathbb{C} \). Suppose \( f \) is an irreducible polynomial in \( F[X,Y] \) with the degree \( \text{deg}_Y f \) of \( f \) in the variable \( Y \) satisfying \( \text{deg}_Y f = n > 0 \). Then

1. For each \( a \) in \( F \), there exists at most \( n \) elements \( b \) in \( F \) satisfying \( (a,b) \in Z_F(f) \).

2. There exists a finite subset \( \Delta \) of \( F \) with the property that for any element \( a \) in \( F \setminus \Delta \), there exist exactly \( n \) elements \( b \) in \( F \) satisfying \( (a,b) \) lies in \( Z_F(f) \).

**Question 41.9.** What is \( Z_F(f) \) if \( f \) in the proposition is irreducible but \( \text{deg}_Y f = 0 \)?

**Proof.** Write \( f = \sum_{i=0}^{n} q_i(X)Y^i \) with \( q_i(X) \) polynomials in \( F[X] \), \( q_n(X) \) nonzero for some \( n > 0 \). If \( a \) lies in \( F \), we have \( f(a,Y) = \sum_{i=0}^{n} q_i(a)Y^i \) has at most \( n \) roots in \( F \) unless \( q_i(a) = 0 \) for all \( i \). But \( X - a \mid q_i(X) \) in \( F[X] \) for all \( i \) if and only if \( X - a \mid f \) in \( F[X,Y] \). As \( f \) is irreducible and \( f \notin F[X] \), this possibility cannot occur. The remaining problem is when \( f(a,Y) \) has a multiple root in \( F[Y] \). Let

\[ \Delta = \{ a \in F \mid q_n(a) = 0 \text{ or } f(a,Y) \text{ has a multiple root in } F \} \]

If \( a \in \Delta \), then \( f(a,Y) \) has at most \( n \) roots in \( F \), as \( f(a,Y) \) is not zero and if \( a \in F \setminus \Delta \) then \( f(a,Y) \) is a product of \( n \) linear polynomials in \( F[Y] \) as \( F \) is algebraically closed, i.e., has precisely \( n \) distinct roots.
So it suffices to show that $\Delta$ is finite. Certainly, there exist finitely many $a$ in $F$ with $q_n(a) = 0$, so it suffices to show the following:

**Claim.** There exist finitely many $a$ in $F$ with the nonzero polynomial $f(a, Y)$ having a multiple root.

As $F$ is algebraically closed, Remark 41.7 says that $f(a, Y)$ has a multiple root if and only if $f(a, Y)$ and $\frac{\partial f(a, Y)}{\partial Y}$ have a common root in $F$.

We know that

(i) $\deg_Y \frac{\partial f(X, Y)}{\partial Y} < \deg_Y f(X, Y)$.

(ii) $f(X, Y)$ is irreducible in $F[X, Y]$ (as $f(X, Y) \in F(X)[Y]$ is irreducible by Lemma 34.7, where $F(X)$ is the quotient field of $F[X]$).

(iii) $\frac{\partial f(X, Y)}{\partial Y}$ is not zero (as $0 < \deg_Y f(X, Y)$ and the characteristic of $F$ is zero (Why?).

These conditions mean that $f$ and $\frac{\partial f}{\partial Y}$ have no common factors in $F[X, Y]$. Consequently, by Proposition 41.3, we have $Z_F(f) \cap Z_F(\frac{\partial f}{\partial Y})$ is a finite set. This proves the claim.

If $F = \mathbb{C}$ in the above, then the affine plane curve $Z_\mathbb{C}(f)$ is called a **Riemann surface**. It is called a surface as $\dim_\mathbb{R} \mathbb{C} = 2$ and it is viewed as a surface over $\mathbb{R}$.

Using a little bit of analysis (e.g., the Implicit Function Theorem), one can show that if $f$ is an irreducible polynomial in $\mathbb{C}[X, Y]$ with $\deg_Y f > 0$, then $Z_\mathbb{C}(f)$ has the following properties: Let $\Delta$ be as in Proposition 41.9 and $Z_\mathbb{C}(f)$ the induced topology in $\mathbb{C}^2$. Then there exists a continuous map $\pi : Z_\mathbb{C}(f) \to Z_\mathbb{C}(Y)$ (so $Z_\mathbb{C}(Y)$ is the $X$-axis in $\mathbb{C}^2$, i.e., the real plane defined by $Y = 0$ in $\mathbb{C}^2$) satisfying

1. $|\pi^{-1}(x)| = n$ for all $x$ in $Z_\mathbb{C}(Y) \setminus \Delta$.
2. For every $x_0$ in $Z_F(Y) \setminus \Delta$, there exists an open neighborhood $U$ of $x$
   such that $\pi^{-1}(U)$ is a union of $n$ disconnected sets $V_i$ (so $\pi^{-1}(U) = V_1 \cup \cdots \cup V_n$) with each $V_i \subset Z_\mathbb{C}(f)$ open and $\pi|_{V_i} : V_i \to U$ a homeomorphism.

We say that $Z_\mathbb{C}(f)$ is an $n$-sheeted branched covering of $Z_F(Y)$.

**Exercises 41.10.**

1. Prove Remark 41.4.
2. Show if $F$ is a field of characteristic zero and $f$ a non-constant polynomial in $F[t]$, the its derivative $f'$ is never zero. Show that if the characteristic of $F$ is not zero, this may not be true by producing counterexamples.
3. Let $F$ be a field of characteristic zero and $f$ a polynomial in $F[t]$. Let $f = (t - a)^r g$ in $F[t]$ for some positive integer $n$ and polynomial $g$ in $F[t]$ with $g(a)$ nonzero. Show that $a$ is a multiple root of $f$, i.e., $r > 1$ if and only if $a$ is a root of $f'$.

4. Show that the Proposition 41.8 still is valid, if we weaken the condition that $f$ be irreducible, to $f$ is irreducible as a polynomial in $X$. 
Finitely Generated Modules Over a PID

This is the deepest, and most difficult chapter in this part of the text. We are interested in classifying finitely generated modules over a PID which we succeed in doing. This has applications to abelian group theory and finite dimension vector spaces.

We begin by discussing row reduction for matrices over a ring. In particular, if the ring is a PID, we shall show that there is a unique matrix up to equivalence of matrices called the Smith Normal Form of the matrix. If $R$ is a euclidean domain, we give an algorithm for computing this matrix in §C. We shall show that any finitely generated module over a PID decomposes into a direct sum of two pieces, a free submodule and a torsion submodule (a module in which for every element $m$ there exists a nonzero ring element $r$ satisfying $rm = 0$). The torsion submodule also decomposes nicely (in two different ways). There is also a uniqueness statement for each of the two forms. For finitely generated abelian groups, this says that, up to isomorphism, every finitely generated abelian group is isomorphic to $\mathbb{Z}^r \times G$ with $r$ unique and $G$ a unique group that is a product of finite cyclic groups, each of which decomposes by the Chinese Reminder Theorem.

The theory is also applied to linear algebra. If $V$ is a finite dimensional vector space over $F$ and $T$ a linear operator on $V$, then $V$ becomes a $F[t]$-module by $t$ acting on $V$ by $T$. This leads to canonical forms of matrices over a field $F$. If $F$ is arbitrary, the desired canonical form is called Rational canonical form. If the field is algebraically closed, the desired canonical form is called Jordan canonical form. Further results in linear algebra are also proven, e.g., the Cayley-Hamilton Theorem.

42. Smith Normal Form

If $R$ is a ring, two $m \times n$ matrices $A$ and $B$ in $R^{m \times n}$ are called equivalent if there exist invertible matrices $P \in \text{GL}_m(R)$ and $Q \in \text{GL}_n(R)$ such that $B = PAQ$. Compare this to Change of Basis Theorem in linear algebra of matrix representations of linear transformations when $R$ is a field (or division ring). [Cf. Appendix B.] Certainly, this relation
is an equivalence relation, so one would like to find a good system of representatives. If \( m = n \), then we are usually interested in the equivalence relation given by similarity of matrices, i.e., when \( Q = P^{-1} \). This last equivalence relation is especially important in linear algebra, i.e., when \( F \) is a field.  [Cf. this to the Change of Basis Theorem for the matrix representation of a linear operator in one basis to another in Appendix B.] For a general field this set of representatives will turn out to be the set of matrices in rational canonical form and over an algebraically closed field this set of representatives will turn out to be the set of Jordan canonical forms. In this section, we determine a set of representatives for equivalence of matrices over a PID. As this includes \( F[t] \) when \( F \) is a field, we shall see that it will aid us solving the equivalence of similarity over matrices. The system of representatives for equivalence of matrices will be the set of matrices in Smith Normal Form. In Appendix C, we will give an algorithm for compute this form of any matrix when \( R \) is a euclidean ring. In this section, we do the more general case that \( R \) is a PID. Unfortunately, unlike the euclidean case, there is no algorithm to compute it.

**Definition 42.1.** Let \( R \) be a domain and \( A \in R^{m \times n} \). We say that \( A \) is in Smith Normal Form if \( A \) is a diagonal matrix of the form

\[
A = \text{diag}(q_1, \ldots, q_r, 0, \ldots, 0) \quad \text{with} \quad q_1 \mid q_2 \mid q_3 \mid \cdots \mid q_r \quad \text{in} \quad R \quad \text{and} \quad q_r \neq 0.
\]

[Note that the number of diagonal entries is the minimum of \( n \) and \( m \).]

We want to prove that every matrix in \( R^{m \times n} \) is equivalent to a matrix in Smith Normal Form if \( R \) is a PID. If \( R \) is a euclidean domain, using the euclidean function provides a method to induct on the value of remainders in the division algorithm. As mentioned above, this leads to a computable algorithm, one used by computer programs to find the Jordan canonical form of square matrices over the complex numbers. We need a substitute for the euclidean function on a euclidean domain. As any PID is a UFD, we can induct on the number of irreducible factors (counted with multiplicity) for any nonzero nonunit.

**Definition 42.2.** Let \( R \) be a UFD and \( a \) a nonzero nonunit in \( R \). Let \( a = p_1 \cdots p_r \) be a factorization into irreducibles in \( R \) (not necessarily distinct). Define the length, \( l(a) \) of \( a \) to be \( r \). If \( a \) is a unit, define its length, \( l(a) \) to be zero.

To prove our result, we need the following lemma. It requires \( R \) to be a PID.
Lemma 42.3. Let \( R \) be a PID and \( A = (a_{ij}) \) a nonzero \( m \times n \) matrix in \( R^{m \times n} \). Fix \( i, j \) with \( a_{ij} \neq 0 \) and \( j < n \). Let \( d \) in the ideal \( (a_{ij}, a_{i,j+1}) \) in \( R \) satisfy \( (d) = (a_{ij}, a_{i,j+1}) \). Then we have:

1. If \( a_{ij} \nmid d \), then \( l(d) < l(a_{ij}) \). In particular, this is true if \( a_{ij} \nmid a_{i,j+1} \).

2. There exists a matrix \( B \) in \( \text{SL}_n(R) \) satisfying \( AB = (c_{ik}) \) with \( c_{ij} = d \) and \( c_{i,j+1} = 0 \).

Remark 42.4. An analogous statement holds for two column entries in \( A \), but in (2), the element \( B \) in \( \text{SL}_n(R) \) multiplying \( A \) on the left.

More generally, a similar statement holds for any two fixed elements in \( A \), at least one not zero, in any fixed row or any fixed column. We prove the lemma as it is stated for convenience of notation.

Proof. (of the lemma). That \( d \) satisfies the first condition is immediate. We construct \( B \) in (2). For notational convenience, write \( v = a_{ij} \) and \( w = a_{i,j+1} \), so \( (d) = (v, w) \). Then \( 0 \neq d = xv + yw = (xa + yb)d \) in the domain \( R \), so \( 1 = xa + yb \) in \( R \), hence

\[
P = \begin{pmatrix} x & -b \\ y & a \end{pmatrix}
\]

lies in \( \text{SL}_2(R) \) [with inverse \( P^{-1} = \begin{pmatrix} a & b \\ -y & x \end{pmatrix} \)].

Then

\[
(v \ w) P = (vx + wy \ -bv + wa) = (d \ -bad + bda) = (d \ 0),
\]

so

\[
B = \begin{pmatrix}
1 & \cdots & 0 \\
\vdots & \ddots & \vdots \\
0 & \cdots & 1
\end{pmatrix}
\]

works (if \( j \neq m \), otherwise have \( P \) end in the \( m \)th row).

\[\square\]

Theorem 42.5. Let \( R \) be a PID and \( A \) an \( m \times n \) matrix in \( R^{m \times n} \). Then \( A \) is equivalent to a matrix in Smith Normal Form.

Proof. We may assume that the matrix \( A = (a_{ij}) \) is nonzero.

We use elementary row and column operations on matrices. We call an elementary row (respectively, column) operation Type I if we add
a multiple of one row (respectively, one column) to another and Type II if we permute two rows (respectively, columns). These correspond to multiplying on the right or left by invertible matrices. [For more details, see Appendix C.] In particular, these elementary row and column operations yield equivalent matrices. Using Type II elementary row and column operations, we may assume that the (1, 1) entry in $A$ is a nonzero element $a$ in $R$ such that among all nonzero entries in $A$, we have $l(a)$ is minimal.

Suppose that $a$ is a unit. Then using Type I elementary row and column operations, we can convert $A$ to the equivalent matrix

$$
\begin{pmatrix}
    a & \cdots & 0 \\
    0 & \ddots & 0 \\
    \vdots & \ddots & A_1 \\
    0 & \cdots & 0
\end{pmatrix}
$$

with $A_1$ an $(m - 1) \times (n - 1)$ matrix.

By induction on matrix size, there exist appropriate invertible matrices $P_1$ and $Q_1$ such that $P_1 A_1 Q_1$ is in Smith Normal Form. Then

$$
\begin{pmatrix}
    1 & \cdots & 0 \\
    0 & \ddots & 0 \\
    \vdots & \ddots & A_1 \\
    0 & \cdots & 0
\end{pmatrix}
\begin{pmatrix}
    a & \cdots & 0 \\
    0 & \ddots & 0 \\
    \vdots & \ddots & Q_1 \\
    0 & \cdots & 0
\end{pmatrix}
\begin{pmatrix}
    1 & \cdots & 0 \\
    \vdots & \ddots & \vdots \\
    \vdots & \cdots & \vdots \\
    \vdots & \cdots & \vdots
\end{pmatrix}
$$

is in Smith Normal Form and is equivalent to $A$. This proves the result in the case that $a$ is a unit. Note that if $A$ is equivalent to a matrix in the form of (*) such that $a$ divides every entry in $A_1$, the argument above shows that we would also be done.

So we may assume that $l(a) > 0$. If $a \nmid a_{ij}$ for some $i = 1$ or $j = 1$, using Lemma 42.3, we can convert $A$ into an equivalent matrix

$$
\begin{pmatrix}
    d & \cdots & * \\
    * & \cdots & * \\
    \vdots & \ddots & \vdots \\
    * & \cdots & *
\end{pmatrix}
$$

with $l(d) < l(a)$, and we are done by induction on $l(a)$. So we may assume that $a \mid a_{ij}$ if $i = 1$ or $j = 1$. Using Type I row and column operations converts $A$ to a matrix equivalent to $A$ of the form (*). By induction of matrix size, there exist invertible matrices of the appropriate size such that $A$ is equivalent to a diagonal matrix $\text{diag}(a, a_2, \ldots, a_r, 0, \ldots, 0)$ with $a_2 \mid a_3 \mid \cdots \mid a_r$ and $a_r \neq 0$ in $R$. If
$a \mid a_2$, we are done; if not use an Type I elementary row operation to convert this matrix to

\[
\begin{pmatrix}
a & a_2 & \cdots \\
0 & a_2 & \\
0 & 0 & a_3 \\
\vdots & \ddots & \\
& & a_r \\
& & \ddots
\end{pmatrix}.
\]

As $a \not\mid a_2$, we can apply the Lemma 42.3 to produce an equivalent matrix with $(1, 1)$ entry $d$ satisfying $l(d) < l(a)$, and we are done by induction on $l(a)$. □

To prove that the Smith Normal Form of a matrix in $R^{m \times n}$ with $R$ a PID is essentially unique, we shall need facts about the determinant that we shall use but shall not prove here. [However, cf. Section 118 for a sophisticated proof of the determinant and its properties.]

Let $R$ be a commutative ring. Then the determinant $\det : \mathbb{M}_n(R) \to R$ is a map that has the following properties:

1. $\det$ is $n$-multilinear (or $n$-linear) as a function of the rows (respectively, columns) of matrices in $\mathbb{M}_n(R)$. This means that it is $R$-linear (i.e., an $R$-homomorphism) in each of the $n$ entries fixing the others, i.e.,

\[
\det(\alpha_1, \ldots, r\alpha_i + \alpha_i', \ldots, \alpha_n) = r \det(\alpha_1, \ldots, \alpha_i, \ldots, \alpha_n) + \det(\alpha_1, \ldots, \alpha_i', \ldots, \alpha_n)
\]

for all $r$ in $R$, and rows (respectively columns) of matrices $A$ in $\mathbb{M}_n(R)$.

2. $\det$ is alternating as a function of the rows (respectively, columns) of matrices in $\mathbb{M}_n(R)$, i.e., if $A$ has two identical rows (respectively, columns), then $\det A = 0$.

[Note: This implies that the matrix obtained by interchanging two rows (respectively columns) of $A$ has determinant $-\det A$.]

3. $\det I = 1$

Indeed it can be shown that $\det : \mathbb{M}_n(R) \to R$ is the unique function satisfying (1), (2), and (3); and it is given by

\[
\det(a_{ij}) = \sum_{\sigma \in S_n} \text{sgn}(\sigma)a_{1\sigma(1)} \cdots a_{n\sigma(n)}
\]
where

$$\text{sgn}(\sigma) = \begin{cases} 
1 & \text{if } \sigma \in A_n \\
-1 & \text{if } \sigma \notin A_n.
\end{cases}$$

Let $A$ be an $m \times n$ matrix in $R^{m \times n}$ and $1 \leq l \leq \min\{m, n\}$. An $l$-order minor of $A$ is a determinant of an $l \times l$ submatrix of $A$, i.e., a matrix obtained from $A$ by deleting $m - l$ rows and $n - l$ columns of $A$. The key to proving our uniqueness statement is the following lemma:

**Lemma 42.6.** Suppose that $R$ is a UFD and $A$ an $m \times n$ matrix in $R^{m \times n}$. Let $P$ and $Q$ be invertible matrices in $\text{GL}_m(R)$ and $\text{GL}_n(R)$ respectively. Set $B = PAQ$. If $1 \leq l \leq \min\{m, n\}$ and

1. if the element $a$ in $R$ is a gcd of all the $l$-order minors of $A$,
and
2. if the element $b$ in $R$ a gcd of all the $l$-order minors of $B$,

then $a$ and $b$ are associates, i.e., the ideals $(a)$ and $(b)$ are the same.

**Proof.** Let $P = (p_{ij})$, $A = (a_{ij})$, $Q = (q_{ij})$ and $c$ in $R$ a gcd of all the $l$-order minors of $PA$. Then the $k$th entry of $PA$ is $\sum_j p_{kj}a_{ji}$, so the $k$th row of $PA$ is $\sum_j p_{kj}(a_{j1}a_{j2}\cdots a_{jn})$ (with the obvious notation). As the determinant is multilinear as a function of the rows of the square submatrices of $A$, we have $a \mid c$ in $R$. As the determinant is multilinear as a function of the columns of the square submatrices of $PA$, an analogous argument shows that $c \mid b$ in $R$, hence $a \mid b$ in $R$. But we also have $A = P^{-1}BQ^{-1}$, so arguing in the same way, we conclude that we also have $b \mid a$. Since $R$ is a domain, $a$ and $b$ must be associates. \qed

**Corollary 42.7.** Suppose that $R$ is a PID and $A$ an $m \times n$ matrix in $R^{m \times n}$ with $B = \text{diag}(d_1, \ldots, d_r, 0, \ldots, 0)$ in $R^{m \times n}$ satisfying $d_1 \mid \cdots \mid d_r$ and $d_r \neq 0$ in $R$ a Smith Normal Form of $A$. Let

$$\Delta_l \text{ be a gcd of all the } l \text{-order minors of } A \text{ in } R \text{ for } 1 \leq l \leq r$$

and $\Delta_0 = 1$. Then

$$\Delta_0 \mid \Delta_1 \mid \cdots \mid \Delta_r \text{ and } d_l \approx \frac{\Delta_l}{\Delta_{l-1}} \text{ in } R \text{ for all } l > 0.$$

Putting this all together, we obtain the following theorem:

**Theorem 42.8.** Let $R$ be a PID and $A$ an $m \times n$ matrix in $R^{m \times n}$. Then $A$ is equivalent to a matrix in Smith Normal Form. Moreover, if $\text{diag}(a_1, \ldots, a_r, 0, \ldots, 0)$ and $\text{diag}(b_1, \ldots, b_s, 0, \ldots, 0)$ are two Smith Normal Forms for $A$, then $r = s$ and $a_i \approx b_i$ for $1 \leq i \leq r$. In particular, the descending sequence of ideals in $R$

$$(a_1) \supset (a_2) \supset \cdots \supset (a_r)$$
completely determine a Smith Normal Form of $A$ with any generator of $(a_l)$ for $1 \leq l \leq r$ being an associate of $\Delta_l/\Delta_{l-1}$, where $\Delta_l$ is a gcd of all the $l$-order minors of $A$ in $R$ and $\Delta_0 = 1$.

The elements $a_1 | \cdots | a_r$ in the theorem are called the invariants factors of $A$. They are unique up to units. If $g : R^m \to R^n$ is an $R$-homomorphism, a matrix representation $[g]_{B,C}$ of $g$ in Smith Normal Form will be called Smith Normal Form of $g$ and the invariant factors of $[g]_{B,C}$ will be called the invariant factors of $g$.

**Exercise 42.9.** Prove Remark 42.4.

### 43. The Fundamental Theorem

In this section, we completely determine finitely generated modules over a PID up to an isomorphism. The basic idea is to first show, unlike in the general case, that any submodule of a finitely generated free module over a PID is itself free. This means that if we have a finitely generated $R$-module $M$ with $R$ a PID, that we have a short exact sequence

$$0 \to R^m \xrightarrow{g} R^n \to M \to 0.$$  

Therefore, $M$ is isomorphic to coker $g$. As $g$ is a map of free $R$-modules, it has a matrix representation $A$ (cf. Appendix B) relative to the appropriate standard bases. Looking at a Smith Normal Form of $A$ produces a direct sum decomposition for coker $g$, hence $M$ decomposes into a direct sum of $R$-cyclic modules. The uniqueness of the Smith Normal Form will show this decomposition is unique up to isomorphism.

We begin by proving that submodules of finitely generated free modules over a PID are free.

**Proposition 43.1.** Let $R$ be a PID, $N$ a free $R$-module of rank $n$, and $M$ a submodule of $N$. Then $M$ is a free $R$-module satisfying $\text{rank } M \leq \text{rank } N = n$.

**Proof.** We know that $N \cong R^n$, so we may assume that $M$ is a submodule of $R^n$; and must show that $M$ is free of rank at most $n$. We show this by induction on $n$. If $n = 0$, this is trivial, so we may assume that $n \geq 1$. Let $\pi_1 : R^n \to R$ given by $(a_1, \ldots, a_n) \mapsto a_1$ be the projection onto the first coordinate, an $R$-epimorphism with kernel $\ker \pi_1 = \{(0, a_2, \ldots, a_n) \mid a_i \in R\} \cong R^{n-1}$ free of rank $n - 1$. Let $\pi_1|_M : M \to R$ be the restriction. It is an $R$-homomorphism with $\ker \pi_1|_M \subset \ker \pi_1 \cong R^{n-1}$, so a free $R$-module
of rank at most \( n - 1 \) by induction. For convenience, let \( \pi'_1 = \pi_1|_M \). In particular, we are done if \( \pi'_1 \) is the zero map. Hence we may assume that \( \pi'_1 \) is not trivial, i.e., \( \emptyset < \text{im} \pi'_1 \subset R \) is a nonzero left ideal in the PID \( R \). It follows that there exists a nonzero \( a \in R \) satisfying \( (a) = \text{im} \pi'_1 \). As \( ra = 0 \) in the domain \( R \) implies that \( r = 0 \), we conclude that \( (a) \cong R \) as \( R \)-modules, so \( (a) \) is free of rank 1. Choose \( m \in M \) such that \( \pi'_1(m) = a \neq 0 \). If \( rm = 0 \), then \( 0 = \pi'_1(rm) = ra \) in the domain \( R \), so \( r = 0 \) and \( Rm \cong R \) is a free \( R \)-module of rank one.

**Claim.** \( M = Rm \oplus \ker \pi'_1 \):

If we show the claim, we are done. Indeed if \( Rm \) is \( R \)-free of rank 1 and \( \ker \pi'_1 \) is \( R \)-free of rank at most \( n - 1 \), then \( M \) is \( R \)-free of rank \( \text{rank} Rm + \text{rank} \ker \pi'_1 \leq n \) by Corollary 38.10, as needed. So we need only show the claim.

We first show \( M = Rm + \ker \pi'_1 \):

Let \( x \in M \), so \( \pi'_1(x) = ra \) for some \( r \in R \), hence lies in \( (a) \). It follows that \( x - rm \) lies in \( \ker \pi'_1 \), hence \( x \) lies in \( Rm + \ker \pi'_1 \).

Next we show \( Rm \cap \ker \pi'_1 = \emptyset \):

Suppose that \( x \) lies in \( Rm \cap \ker \pi'_1 \). Then \( x = rm \) for some \( r \in R \) and \( 0 = \pi'_1(x) = \pi'_1(rm) = ra \) in the domain \( R \). As \( a \) is nonzero, we have \( r = 0 \), so \( x = 0 \) as needed.

This establishes the claim and hence the proposition. \( \square \)

As mentioned before, it is in fact true that any submodule of a free module over a PID is free, finitely generated or not. The fact that a free submodule of a finitely generated free module \( N \) has rank bounded by the rank of \( N \) is in fact true over any commutative ring, but this is harder (especially the case when the ring is not a domain).

The theorem has an interesting consequence even stronger than that mentioned in the introduction of this section.

**Corollary 43.2.** Let \( R \) be a PID and \( M \) a finitely generated \( R \) module. Then there exists an exact sequence of \( R \)-modules

\[
0 \to R^m \xrightarrow{g} R^n \xrightarrow{f} M \to 0. 
\]

with \( m \leq n \).

**Proof.** As \( M \) is generated by \( n \) elements, for some positive integer \( n \), we have an exact sequence

\[
0 \to \ker f \xrightarrow{\text{inc}} R^n \xrightarrow{f} M \to 0. 
\]

As \( \ker f \subset R^n \), it is \( R \)-free of rank at most \( n \) by the proposition. \( \square \)
If $M$ is a finitely generated $R$ module with $R$ a PID, then the sequence

$$0 \rightarrow R^m \xrightarrow{g} R^n \xrightarrow{f} M \rightarrow 0.$$ 

says that $M$ can be generated by $n$ elements and the relations on these generators generates a free $R$-module of rank at most $m$, i.e., the minimal number of relations on the generators that generate all relations on these generators is at most $m$. It is called a free resolution of $M$. We shall also say that $M$ has a free presentation. In particular, if $n$ is the minimal number of generators of $M$, one can show that $m$ is the minimal number of generating relations for $M$ among all free resolutions of $M$. We can also say something about submodules of finitely generated modules over a PID.

**Corollary 43.3.** Let $R$ be a PID and $M$ a finitely generated $R$-module. Suppose that $M$ can be generated by $n$ elements. Then any submodule of $M$ can be generated by $n$ elements.

**Proof.** Let

$$0 \rightarrow R^m \xrightarrow{g} R^n \xrightarrow{f} M \rightarrow 0$$

be a free resolution of $M$ and $N$ a submodule of $M$. By the Correspondence Principle, there exists a submodule $B$ of $R^n$ satisfying $\ker f \subset B \subset R^n$ and $f(B) = N$, i.e., $N \cong B/\ker f$. As $B$ is $R$-free of rank at most $n$, it can be generated by $\leq n$ elements, hence the same is true of $N$. \qed

**Remarks 43.4.** 1. If $R$ is a commutative noetherian ring but not a PID, then there exists a non-principal ideal $\mathfrak{A} = (x,y)$ of $R$. We have an exact sequence

$$0 \rightarrow \mathfrak{A} \rightarrow R \xrightarrow{\pi} R/\mathfrak{A} \rightarrow 0.$$ 

$R$ is a cyclic $R$-module but $\mathfrak{A}$ is not. Nor is $\mathfrak{A}$ $R$-free as $x \cdot y - y \cdot x = 0$ in $R$ for any generators of $\mathfrak{A}$. In fact, a similar argument shows that an ideal in a ring $R$ is free as an $R$-module if and only if it is principal. This shows that the last three results do not hold for commutative rings that are not PIDs.

2. The proper $\mathbb{Z}$-free submodule $2\mathbb{Z}$ of $\mathbb{Z}$ is of rank 1, the same as $\mathbb{Z}$.

3. It is, in fact, true that if $R$ is a commutative ring, $M$ a finitely generated free $R$-module of rank $n$, then any submodule of $M$ that is also $R$-free has rank at most $n$. The proof is not easy (although easier for the case when $R$ is a domain).
Observation 43.5. Let \( N_i \) be a submodule of the \( R \)-module \( M_i \) for \( i = 1, 2 \). Then we have an \( R \)-isomorphism
\[
(M_1 \coprod M_2)/(N_1 \coprod N_2) \cong (M_1/N_1) \coprod (M_2/N_2)
\]

Theorem 43.6. (Fundamental Theorem of fg Modules over a PID, Form I) Let \( R \) be a PID and \( M \) a nontrivial finitely generated \( R \)-module. Then \( M \) is a direct sum of cyclic \( R \)-modules. More precisely, there exist a nonnegative integer \( r \) and
\[
(*) \quad m_1, \ldots, m_r \in M \text{ satisfying } M = P \oplus Rm_1 \oplus \cdots \oplus Rm_r
\]
with \( P \) a free \( R \)-module of rank \( s \), some \( s \geq 0 \), and there exist nonzero nonunits
\[
(**) \quad d_i \in R \text{ satisfying } \begin{cases}
\text{ann}_R m_i = (d_i) & i = 1, \ldots, r \\
d_1 | \cdots | d_r.
\end{cases}
\]
Moreover, \( r \) and \( s \) are unique in (*) and the \( d_i \) in (**) are unique up to units and the order determined by \( d_1 | \cdots | d_r \), i.e., any direct sum decomposition of \( M \) satisfying (*) and (**) has the \( R \)-free submodule in (*) of the same rank \( s \), the number of non-free \( R \)-cyclic modules the same \( r \) with the corresponding descending chain of ideals
\[
(d_1) \supset (d_2) \supset \cdots \supset (d_r)
\]
unique.

The elements \( d_1, \ldots, d_r \) in the Fundamental Theorem 43.6 are called invariant factors of \( M \). Note also that the \( Rm_i \) in the Fundamental Theorem 43.6 are unique up to isomorphism with \( Rm_i \cong R/(d_i) \). Of course, if \( d_i \approx d_j \), then \( Rm_i \cong Rm_j \). In particular,
\[
M \cong R^s \coprod R/(d_1) \coprod \cdots \coprod R/(d_r),
\]
and the right hand side completely determines the isomorphism class of \( M \), i.e., the invariant factors and \( s \) completely determine \( M \) up to isomorphism.

**Proof. Existence:** We have shown that there exists an exact sequence
\[
(43.7) \quad 0 \to R^n \xrightarrow{g} R^m \to M \to 0
\]
of \( R \)-modules with \( m \leq n \), i.e., \( M \cong \text{coker } g = R^n/\text{im } g \). Let \( S_m \) and \( S_n \) be the standard bases for \( R^m \) and \( R^n \), respectively. By Theorem 42.8 and the Change of Basis Theorem from linear algebra (cf. Theorem B.7 in Appendix B), there exist bases \( \mathcal{B} \) and \( \mathcal{C} \) for \( R^m \) and \( R^n \), respectively, such that the \( n \times m \) matrix \([g]_{\mathcal{B}, \mathcal{C}}\) is a Smith Normal Form of \([g]_{S_m, S_n}\).
Let $C = \{v_1, \ldots, v_n\}$ and $[g]_{B,C} = \text{diag}(d_1, \ldots, d_r, 0, \ldots, 0)$ with $d_1 | \cdots | d_r$ in $R$, $r \leq m \leq n$, and $d_r$ nonzero. As $g$ is injective, no column of $[g]_{B,C}$ can be zero. Therefore, $r = m$. We have 

$$\text{im } g = \bigoplus_{i=1}^{m} Rd_i v_i = \bigoplus_{i=1}^{m} Rd_i v_i \bigoplus \bigoplus_{i=m+1}^{n} 0.$$ 

let $d_i = 0$ for $m < i \leq n$. As $R \to R v_i$ by $r \mapsto rv_i$ is an $R$-isomorphism, we have $R/(d_i) \cong R v_i/Rd_i v_i$ for $i = 1, \ldots, n$. Hence by Observation 43.5, we have 

$$M \cong \text{coker } g \cong (\prod_{i=1}^{n} R v_i)/(\prod_{i=1}^{n} Rd_i v_i) \cong \prod_{i=1}^{m} (R v_i/Rd_i v_i) \prod_{i=m+1}^{n} R v_i \cong \prod_{i=1}^{m} R/(d_i) \prod_{i=1}^{n} R^{n-m}.$$ 

If $d_i$ is a unit, then $R/(d_i) = 0$, so dropping such $d_i$, if 

$$f : M \longrightarrow \prod_{d_i \neq R^\times}^{m} \frac{R}{(d_i)} \prod_{i=1}^{n} R^{n-m}$$ 

is the isomorphism, then $f^{-1}$ gives the desired decomposition with $m_i = f^{-1}(1_R + (d_i))$ for $d_i$ not a unit or zero and $P = f^{-1}(R^{n-m})$.

For the uniqueness, we need further observations that we leave as exercises.

**Observations 43.8.** Let $R$ be a PID, $e$ and $d$ nonzero elements in $R$, $g$ a gcd of $e$ and $d$, $p$ a prime element in $R$ (so $p$ is nonzero and $R/(p)$ is a field), and $N$ an $R$-module. Then the following are true:

1. $dR^m \cong R^m$.
2. $d(R/(e)) \cong R/(e)$.
3. $R/(p)$ acts on $N/pN$ by 
   $$(r + (p))(x + pN) := rx + pN \quad = r(x + pN)$$ 
   for all $r \in R$ and $x \in N$, i.e., the $R/(p)$-action and the $R$-action on $N/pN$ are compatible (have the same effect). In particular, $N/pN$ is a vector space over $R/(p)$.
4. As $R$-modules and $R/(p)$-vector spaces, we have 
   $$(R/(d))/p(R/(d)) \cong \begin{cases} R/(p) & \text{if } p \mid d \\ 0 & \text{if } p \nmid d. \end{cases}$$
Using these observations, we show:

**Uniqueness:** Suppose that
\[ R/(d_1) \prod R/(d_r) \prod R^s \cong M \cong R/(d'_1) \prod R/(d'_{r'}) \prod R^{s'} \]
with \( d_1 | \cdots | d_r \), the \( d_i \) nonzero nonunits and \( d'_1 | \cdots | d'_{r'} \), the \( d'_i \), nonzero nonunits, for some integers \( r, r', s, s' \).

We know that a submodule of a free \( R \)-module is \( R \)-free, hence an isomorphic copy of \( R/(d) \) with \( d \) a nonzero nonunit in \( R \) cannot be a submodule of an \( R \)-free. In particular, if there are no \( d_i \) or no \( d'_i \), then \( M \) is free and there are neither any \( d_i \) nor any \( d'_i \). As \( R \) is commutative, \( s = \text{rank } M = s' \) by Theorem 38.9. So we may assume both \( r \geq 1 \) and \( r' \geq 1 \). As \( d_r (R/(d_i)) = 0 \) and \( d'_{r'} (R/(d'_i)) = 0 \) for all \( i \) and \( i' \). By Observation 43.8(1), we have
\[ R^s \cong d_r d'_r, R^s \cong d_r d'_r, M \cong d_r d'_r R^{s'} \cong R^{s'} \]
is \( R \)-free, hence \( s = s' \) by Theorem 38.9. Let \( m \in M \) correspond to a generator of \( R/(d'_r) \) under the given isomorphism. Then \( Rd_r m \subset d_r M \cong R^s \) must also be \( R \)-free. As \( d'_r d_r m = 0 \), we have \( d_r m = 0 \) (Why?), so \( d_r \in \text{ann}_R m = (d'_r) \). Similarly, \( d'_r \in (d_r) \), so \( d_r \cong d'_r \), i.e., \( (d_r) = (d'_r) \).

Assume that \( r \geq r' > 0 \). Let \( p \) be a prime element (so nonzero) satisfying \( p \mid d_1 \), which exists as \( d_1 \) is a nonzero nonunit. Then by Observations 43.8(2) and (4) and Observation 43.5, we have \( M/pM \cong \prod_{i=1}^{r+s} R/(p) \) is an \( R/(p) \)-vector space, hence of rank \( r + s \). Therefore, \( \prod_{i=1}^{r'} R/(d'_i)/p(R/(d'_i)) \prod R^s/pR^s \) is free of rank \( r + s \). But, by Observation 43.8(4), we also know that it is free of rank at most \( r' + s \). It follows that \( r = r' \). In particular, we are done if \( p \approx d_r \approx d'_r \). By Observations 43.8(1) and (2), we have
\[ \prod R/(d_i/p) \prod R^s \cong pM \cong \prod R/(d'_i/p) \prod R^s. \]
By induction on the length \( l(d_r) \) of \( d_r \), we conclude that \( \frac{d_i}{p} \approx \frac{d'_i}{p} \) for all \( i \), hence \( d_i \approx d'_i \) for all \( i \). \( \square \)

Note the proof of the uniqueness statement is similar to the corresponding proof of the Fundamental Theorem of Arithmetic.

**Remark 43.9.** It should also be noted, that to prove the existence in Theorem 43.6, we do not need to use the existence of a free presentation of \( M \). Indeed the same proof works if we have an exact sequence
\[ R^m \xrightarrow{f} R^n \rightarrow M \rightarrow 0. \]
But such an exact sequence always exists by Corollary 39.8, since a PID is a Noetherian ring.

**Remarks 43.10.** Let $M$ be a finitely generated module over a PID.

1. The proof of the uniqueness part of Theorem 43.6 shows, up to isomorphism, $M$ depends only on the cokernel of $g$ in the free presentation (43.7) of $M$. Also note that $\text{im } g \cong R^m$ is of rank $m$. This determination of $M$ up to isomorphism is independent of the integers $m$ and $n$. Indeed given any other free presentation

$$0 \to R^{m'} \xrightarrow{g'} R^{n'} \to M \to 0,$$

we know by Proposition 43.1, that we also must have $m' \leq n'$. Moreover, if we compute the Smith Normal form associated to the map $g'$, we must have $\text{coker } g \cong \text{coker } g'$ by the uniqueness argument in the proof of Theorem 43.6, thus determining $M$ up to isomorphism, for this different $n'$. The only difference would be in the number of $d_i$ that are units arising in each of the two Smith Normal Forms arising from $g$ and $g'$, respectively, i.e., $|n - m|$ and $|n' - m'|$, respectively; and we are dropping the trivial modules corresponding to such $d_i$.

2. The only remaining question is what if we have an exact sequence

$$R^m \xrightarrow{g'} R^n \to M \to 0$$

with $g'$ not necessarily an $R$-monomorphism? In this case, we have exact sequences

$$0 \to \ker g' \to R^m \xrightarrow{g'} R^n \to M \to 0.$$

and

$$0 \to R^m / \ker g' \xrightarrow{\overline{g'}} R^n \to M \to 0$$

with $\overline{g'}$ the canonical epimorphism induced by First Isomorphism Theorem. The second exact sequence is a free presentation since $R^m / \ker g'$ is free of rank at most $n$ and $\text{coker } g' = \text{coker } \overline{g'}$. Indeed, the only difference arising between these two will be zero columns in the Smith Normal form of $g'$ arising from its kernel.

**Examples 43.11.** 1. We determine the abelian group $M$, up to isomorphism, generated by $x_1, x_2, x_3, x_4$ subject to the relations

$$2x_1 + 3x_2 + x_4 = 0$$

$$x_1 - 2x_2 + x_3 = 0.$$
It is easy to see that the following is an exact sequence:

\[ \mathbb{Z}^2 \xrightarrow{g} \mathbb{Z}^4 \xrightarrow{h} M \to 0 \]

with \( g = \begin{pmatrix} 2 & 1 \\ 3 & -2 \\ 0 & 1 \\ 1 & 0 \end{pmatrix} \) and \( h \) takes \( e_i \mapsto x_i, \ 1 \leq i \leq 4 \), i.e., \( M \cong \text{coker } g \). To find the isomorphism type of \( M \) we find a Smith Normal Form of \( g \). Using the algorithm in Appendix C (essentially the division algorithm), we row and column reduce \( g \) to

\[ \begin{pmatrix} 1 & 0 \\ 0 & 1 \\ 0 & 0 \\ 0 & 0 \end{pmatrix} \]

Hence by Remark 43.10(1), \( M \) is isomorphic to \( \mathbb{Z}/\mathbb{Z} \uplus \mathbb{Z}/\mathbb{Z} \uplus \mathbb{Z}/\mathbb{Z} \uplus \mathbb{Z} \cong \mathbb{Z}^2 \). We also conclude that

\[ 0 \to \mathbb{Z}^2 \xrightarrow{g} \mathbb{Z}^4 \xrightarrow{h} M \to 0 \]

is a free resolution of \( M \).

2. We determine the abelian group \( M \), up to isomorphism, generated by \( x_1, x_2, x_3, x_4 \) subject to the relations

\begin{align*}
4x_1 + 2x_2 + 4x_3 + 3x_4 &= 0 \\
2x_1 + 2x_2 - 2x_3 + 2x_4 &= 0 \\
-6x_1 + 6x_3 - 6x_4 &= 0
\end{align*}

We have an exact sequence

\[ \mathbb{Z}^3 \xrightarrow{g} \mathbb{Z}^4 \xrightarrow{h} M \to 0 \]

with \( g = \begin{pmatrix} 4 & 2 & -6 \\ 2 & 2 & 0 \\ 4 & -2 & 6 \\ 3 & 2 & -6 \end{pmatrix} \) and \( h \) takes \( e_i \mapsto x_i, \ 1 \leq i \leq 4 \), Row and column reducing \( g \) using the division algorithm as in Appendix C, we see that a Smith Normal Form for \( g \) is

\[ \begin{pmatrix} 1 & 0 & 0 \\ 0 & 2 & 0 \\ 0 & 0 & 6 \\ 0 & 0 & 0 \end{pmatrix} \] and \( h \) takes \( e_i \mapsto x_i, \ 1 \leq i \leq 4 \),
Hence, by Remark 43.10(1), $M$ is isomorphic to
\[ \mathbb{Z}/\mathbb{Z} \bigoplus \mathbb{Z}/2\mathbb{Z} \bigoplus \mathbb{Z}/6\mathbb{Z} \cong \mathbb{Z}/2\mathbb{Z} \bigoplus \mathbb{Z}/6\mathbb{Z}. \]
Note that this is also isomorphic to $\mathbb{Z}/2\mathbb{Z} \bigoplus \mathbb{Z}/2\mathbb{Z} \bigoplus \mathbb{Z}/3\mathbb{Z} \bigoplus \mathbb{Z}$. This last isomorphism is what the second form of the Fundamental Theorem will give. We also conclude that
\[ 0 \to \mathbb{Z}^3 \xrightarrow{g} \mathbb{Z}^4 \xrightarrow{h} M \to 0 \]
is a free resolution of $M$.

We next look at a property of modules that is, in general, weaker than being free but for finitely generated modules over a PID turn out to be equivalent.

**Definition 43.12.** Let $R$ be a domain an $M$ and $R$-module. An element $m$ in $M$ is called an $R$-torsion element if there exists a nonzero element $r$ in $R$ satisfying $rm = 0$, i.e., $\text{ann}_R m > 0$. Set
\[ M_t := \{ m \in M \mid m \text{ is an } R\text{-torsion element} \}. \]
We say that $M$ is a torsion $R$-module if $M = M_t$ and a torsion-free $R$-module if $M_t = 0$.

The proof of the following is straight-forward and left as an exercise.

**Properties 43.13.** Let $R$ be a domain, $M$ an $R$-module, and $m$ a nonzero element in $M$. Then
1. The element $m$ is not an $R$-torsion element if and only if $Rm$ is torsion-free if and only if $Rm$ is $R$-free.
2. $M_t$ is a submodule of $M$.
3. $M/M_t$ is a torsion-free $R$-module.
4. If $M$ is $R$-free, then it is $R$-torsion-free.

**Remarks 43.14.** 1. Any finite abelian group is a $\mathbb{Z}$-torsion module.
2. If $S$ is a domain and $R \subset S$ is a subring, then $S$ is an $R$-torsion-free module.
3. $\mathbb{Q}$ is $\mathbb{Z}$-torsion-free but not $\mathbb{Z}$-free.

We shall show that over a PID finitely generated torsion-free modules are free. The following computation together with the Fundamental Theorem 43.6 will ensure this.

**Example 43.15.** Let $R$ be a PID and $d_1 | \cdots | d_r$ nonzero nonunits in $R$. Set
\[ N = R/(d_1) \bigoplus \cdots \bigoplus R/(d_r) \bigoplus R^* \] and
\[ N_0 = R/(d_1) \bigoplus \cdots \bigoplus R/(d_r), \]
so \( N = N_0 \sqcup R^s \). We have
\[
d_r N = d_r N_0 \sqcup d_r R^s = d_r R^s \cong R^s
\]
is \( R \)-free of rank \( s \) and \( N_0 \subset N_t \). Next, suppose that \((n_0, v) \in N_t\) with \(n_0 \in N_0\) and \(v \in R^s\). Then there exists a nonzero element \( a \) in \( R \) satisfying \( a(n_0, v) = 0 \), i.e., \( an_0 = 0 \) and \( av = 0 \). As \( R^s \) is \( R \)-free, it is \( R \)-torsion-free, so \( v = 0 \) and we conclude that \( N_0 = N_t \). As the diagram
\[
\begin{array}{c}
0 \longrightarrow N_0 \longrightarrow N \longrightarrow R^s \longrightarrow 0 \\
\end{array}
\]
has exact rows and is commutative, the vertical right hand arrow exists, (why?), hence is an isomorphism, so \( N/N_t \cong R^s \). In particular, \( N \) is \( R \)-free if and only if \( N_0 = 0 \) if and only if \( N_t = 0 \) if and only if \( N \) is \( R \)-torsion-free.

The example together with the Fundamental Theorem yields:

**Corollary 43.16.** Let \( R \) be a PID and \( M \) a finitely generated \( R \)-module. Then

1. \( M \) is \( R \)-torsion-free if and only if \( M \) is \( R \)-free.
2. \( M = M_t \oplus P \) for some finitely generated free \( R \)-module \( P \subset M \) with rank \( P \) unique.

**Remark 43.17.** \( \mathbb{Q} \) is a torsion-free abelian group but not \( \mathbb{Z} \)-free, so the assumption in the corollary that the module be finitely generated torsion-free is necessary in general to conclude that it is free. [The abelian group \( \bigoplus_{i=1}^{\infty} \mathbb{Z} \) is another example of a torsion-free abelian group that is not \( \mathbb{Z} \)-free.]

We turn to establishing an alternate form of the Fundamental Theorem. We must do some preliminary work. We begin with another example.

**Example 43.18.** Let \( R \) be a PID and \( M \) a nontrivial cyclic \( R \)-module. Then \( M = Rm \cong R/(d) \), for some \( m \in M \) and unique \( (d) = \text{ann}_R m < R \). Assume that \( M \) is not \( R \)-free. Then \((d) > 0 \) and \( d \) is not a unit. Let \( d = p_1^{e_1} \cdots p_r^{e_r} \) be a factorization, with \( p_1, \ldots, p_r \) non-associative irreducible, hence prime, elements and \( e_1, \ldots, e_r \) positive integers. By the Chinese Remainder Theorem, \( R/(d) \cong R/(p_1^{e_1}) \times \cdots \times R/(p_r^{e_r}) \) as rings. It follows that \( R/(d) \cong R/(p_1^{e_1}) \sqcup \cdots \sqcup R/(p_r^{e_r}) \) as \( R \)-modules, as the \( R \)-action on \( R/(a) \) is compatible with the \( R/(a) \)-action for any \( a \in R \). (Cf. Observation 43.8(3).) Therefore, there exist \( m_1, \ldots, m_r \).
in $M$ satisfying $M = Rm_1 \oplus \cdots \oplus Rm_r$ with $\text{ann}_R m_i = (p_i^{e_i})$ for $i = 1, \ldots, r$.

Although this example suffices to get our alternate form, we first make some further comments. If $R$ is a domain, $p$ an irreducible element in $R$, and $M$ an $R$-module, we say that $M$ is $p$-primary if for all $x$ in $M$, there exists a positive integer $n$ (depending on $x$) satisfying $p^n x = 0$. If $R$ is Noetherian, let $\mathcal{P}$ denote a system of representatives of equivalence classes of irreducible elements in $R$ under the equivalence relation $\sim$.

**Example 43.19.** Let $R$ be a Noetherian domain and $p \in \mathcal{P}$. So $R/(p^r)$ is a $p$-primary $R$-module for all positive integers $r$.

**Theorem 43.20.** (Primary Decomposition Theorem) Let $R$ be a PID and $M$ a torsion $R$-module. If $p \in \mathcal{P}$, let

$$M_p := \{x \in M \mid p^r x = 0 \text{ for some positive integer } r\}.$$ 

Then $M_p$ is a $p$-primary submodule of $M$ and $M = \bigoplus_{p \in \mathcal{P}} M_p$. Moreover, if $M$ is finitely generated, then $M_p = 0$ for almost all $p \in \mathcal{P}$.

**Proof.** Clearly, $M_p$ is a $p$-primary submodule of $M$ for all $p \in \mathcal{P}$. We may assume that $M$ is not trivial.

$M = \sum_{p \in \mathcal{P}} M_p$: Let $x$ be a nonzero element of $M$ and $(d) = \text{ann}_R x$. Then $d$ is nonzero nonunit as $M$ is a torsion $R$-module. It follows by Example 43.18 that $x$ lies in $\sum_{p \in \mathcal{P}} M_p$.

$M = \bigoplus_{p \in \mathcal{P}} M_p$: Suppose that $x \in M_{p_0} \cap \sum_{p \neq p_0 \in \mathcal{P}} M_p$. Let $(d) = \text{ann}_R x$, so $M_p = 0$ for all $p \in \mathcal{P}$, and positive integers $e_1, \ldots, e_n$ satisfying $x = x_1 + \cdots + x_n$ with $p_1^{e_1} x_i = 0$ for $i = 1, \ldots, n$. It follows that $p_1^{e_1} \cdots p_n^{e_n} x = 0$, so $p_1^{e_1} \cdots p_n^{e_n} \in \text{ann}_R x$, i.e., $d \mid p_1^{e_1} \cdots p_n^{e_n}$. By choice $p_0$ and $p_1^{e_1} \cdots p_n^{e_n}$ are relatively prime in the PID $R$. It follows that $d$ is a unit hence $x = 0$.

Finally, suppose that $M$ is finitely generated, say $M = Rx_1 + \cdots + Rx_n$ with $x_i$ nonzero for $i = 1, \ldots, n$. Let $d_i$ be a nonzero nonunit satisfying $(d_i) = \text{ann}_R x_i$ for $i = 1, \ldots, n$ and set $d = d_1 \cdots d_n$ a nonzero element in the domain $R$. It follows that $M_p = 0$ for all $p \in \mathcal{P}$ satisfying $p \not{|} d$. Hence $M = \bigoplus_{p \in \mathcal{P}} M_p$, a finite sum. \qed

**Example 43.21.** Suppose that $R$ is a PID and $M$ is a nontrivial cyclic $R$-module with $\text{ann}_R x = (d)$ and $d = p_1^{e_1} \cdots p_r^{e_n}$, with $p_i \in \mathcal{P}$. Then $M = \bigoplus_{i=1}^n M_{p_i}$ and $M_{p_i} \cong R/(p_i^{e_i})$ for $i = 1, \ldots, n$. 

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Using this, together with the Fundamental Theorem 43.6, we can reformulate it as

**Theorem 43.22.** (Fundamental Theorem of fg Modules over a PID, Form II) Let $R$ be a PID and $M$ a nontrivial finitely generated $R$-module. Then there exists a finitely generated free submodule $P$ of $M$ of unique rank and $p_i \in \mathcal{P}$, $i = 1, \ldots, r$ some $r$ and $x_{ij}$ in $M$, $1 \leq j \leq n_i$ and $1 \leq i \leq r$ and for some $n_i$, satisfying

$$M = P \oplus \left( \bigoplus_{i=1}^{r} \bigoplus_{j=1}^{n_i} Rm_{ij} \right)$$

and for all $i$ and $j$, we have

1. $Rm_{ij} \cong R/(p_i^{e_{ij}})$ and
2. $1 \leq e_{i1} \leq e_{i2} \leq \cdots \leq e_{in_i}$.

Moreover, this decomposition is unique relative the ordering of the $p_i$ subject to the order of the $e_{ij}$ above and the uniqueness of the rank of $P$.

Generators of the ideals $(p_i^{e_{ij}})$ in the theorem are called elementary divisors of $M$. They are unique up to units, and together with the rank of $P$ completely determine $M$ up to isomorphism. The condition on the $e_{ij}$, of course, orders the set of cyclic modules that are $p_i$-primary submodules in the decomposition (which are unique up to isomorphism). This data completely determines the isomorphism class of $M$.

Of course, both forms of the Fundamental Theorem apply to finitely generated abelian groups, although in this case an easier proof will suffice. If we are given any finitely generated abelian group defined by generators and relations, the invariant factors will determine completely the group up to isomorphism. Moreover, given a positive integer $n$, up to isomorphism, we can write down every abelian group of order $n$ (assuming that we can factor $n$). We give an example how this works.

**Example 43.23.** Find all abelian groups up to isomorphism of order $2^3 \cdot 3^2 \cdot 5$. We write these groups in both forms. Given the factorization, it is easier to write down the representatives of isomorphism classes in Form II. To get Form I, we take the biggest $p$-primary pieces for each prime $p$ and combine them to get one term. Then continue with what is left. This is illustrated in our example by the overlines and underlines of various $p$-primary pieces and what they become in Form I.
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\[ \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/3\mathbb{Z} \times \mathbb{Z}/3\mathbb{Z} \times \mathbb{Z}/5\mathbb{Z} \cong \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/6\mathbb{Z} \times \mathbb{Z}/30\mathbb{Z} \]

\[ \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/3\mathbb{Z} \times \mathbb{Z}/9\mathbb{Z} \times \mathbb{Z}/5\mathbb{Z} \cong \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/90\mathbb{Z} \]

\[ \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/4\mathbb{Z} \times \mathbb{Z}/3\mathbb{Z} \times \mathbb{Z}/5\mathbb{Z} \cong \mathbb{Z}/6\mathbb{Z} \times \mathbb{Z}/60\mathbb{Z} \]

\[ \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/9\mathbb{Z} \times \mathbb{Z}/3\mathbb{Z} \times \mathbb{Z}/5\mathbb{Z} \cong \mathbb{Z}/3\mathbb{Z} \times \mathbb{Z}/120\mathbb{Z} \]

\[ \mathbb{Z}/8\mathbb{Z} \times \mathbb{Z}/9\mathbb{Z} \times \mathbb{Z}/5\mathbb{Z} \cong \mathbb{Z}/360\mathbb{Z}. \]

[Note: Abelian direct sums of abelian groups are usually written as products as a finite product of abelian groups is the same as a direct sum of the same groups.]

**Exercises 43.24.**

1. Let \( R \) be a PID. Suppose that \( M \) is a finitely generated \( R \)-module and

\[ 0 \to R^m \xrightarrow{g} R^n \to M \to 0 \]

is a free presentation of \( M \) with \( n \) minimal. Without using the Fundamental Theorem (or Smith Normal Forms) show if

\[ 0 \to R^{m'} \xrightarrow{g'} R^{n'} \to M \to 0 \]

is another free presentation of \( M \), then \( m' \leq m \).

2. Prove the following generalization of Proposition 43.1 using Exercises 38.12(11), (13)): Let \( R \) be a commutative (this is not needed) ring with the property that every (left) ideal of \( R \) is \( R \)-projective. Let \( M \) be an \( R \)-module isomorphic to a submodule of \( R^n \). Then \( M \) is a direct sum of \( m \leq n \) submodules each isomorphic to a (left) ideal of \( R \). In particular \( M \) is \( R \)-projective.

3. Prove Observation 43.5.

4. Prove all the statements in Observations 43.8.

5. Prove all the properties in 43.13.

6. Let \( R \) be a domain. Show every ideal in \( R \) is \( R \)-torsion-free, but is free if and only if it is principal. In particular, show that \( R \) is a PID if every submodule of a free module is free. [The converse is also true.]

7. Let \( R \) be a domain and \( M \) a finitely generated \( R \)-module. Prove that \( M \) is isomorphic to a free submodule of finite rank.
8. Let \( f : \mathbb{Z}^n \to \mathbb{Z}^m \) be a \( \mathbb{Z} \)-module homomorphism. Let \( S_i \) be the standard basis for \( \mathbb{Z}^i \). Prove that \( f \) is monic if and only if the rank of \([f]_{S^n, S^m}\) is \( n \) and \( f \) is epic if and only if a gcd of the \( m \)th ordered minors of \([f]_{S^n, S^m}\) is 1.

9. Let \( R \) be a commutative ring. Let \( E_n(R) \) be the subgroup of \( \text{GL}_n(R) \) generated by all matrices of the form \( I + \lambda \) where \( \lambda \) is a matrix with precisely one non-zero entry and this entry does not occur on the diagonal. Suppose that \( R \) is a Euclidean ring. Show that \( \text{SL}_n(R) = E_n(R) \). (Cf. Appendix C.)

10. Let \( A \) be a finite abelian group and let

\[ \hat{A} := \{ \chi : A \to \mathbb{C}^\times \mid \chi \text{ a group homomorphism} \}. \]

It is easily checked that \( \hat{A} \) is a group via \( \chi_1\chi_2(x) := \chi_1(x)\chi_2(x) \). Show

(i) \( A \) and \( \hat{A} \) have the same order and, in fact, are isomorphic.

(ii) If \( \chi \) is not the identity element of \( \hat{A} \) then \( \sum_{a \in A} \chi(a) = 0 \).

11. Determine all abelian groups of order 400 up to isomorphism.

### 44. Canonical Forms for Matrices

In this section, we use the Fundamental Theorem of fg Modules over a PID to determine a good system of representatives for matrices over a field under the equivalence relation of similarity. Recall if \( V \) is a finite dimensional vector space over a field \( F \) with basis \( B \) and \( T : V \to V \) is a linear operator then the characteristic polynomial of \( T \) is defined to be

\[ f_T := \det(tI - [T]_B) \in F[t]. \]

This is independent of bases by the Change of Basis Theorem (cf. Appendix B). Roots of \( f_T \) are called eigenvalues of \( T \). An element \( \lambda \) is a root of \( f_T \) if and only if there exists a nonzero vector \( v \) in \( V \) such that \( T(v) = \lambda v \). Such a \( v \) is called a (nonzero) eigenvector for \( T \).

**Construction 44.1.** Let \( V \) be a finite dimensional vector space of dimension \( n \) over a field \( F \) and \( T : V \to V \) a linear operator (endomorphism). Define

\[ \varphi : F[t] \to \text{End}_F(V) \text{ be given by } f \mapsto f(T) \]

evaluation at \( T \). This is a ring homomorphism with

\[ \text{im } \varphi = \{ \sum_{i=0}^{r} a_i T^i \mid a_i \in F \text{ some } m \}, \]
a commutative subring of \( \text{End}_F(V) \). By the First Isomorphism Theorem \( F[t]/\ker \varphi \cong \text{im } \varphi = F[T] \). As \( F[t] \) is a vector space over \( F \) on basis \( \{ t^i \mid i \geq 0 \} \), we have \( \dim_F F[t] \) is infinite. We also know \( \dim_F F[T] \leq \dim_F \text{End}_F(V) = \dim_F M_n(F) = n^2 \) is finite. Thus \( \varphi \) cannot be monic. Consequently, \( 0 < \ker \varphi \). As \( \varphi \) is a ring homomorphism, \( \varphi(1) = 1_V \), so \( \varphi \) is not the trivial map. Therefore, \( F[t] \) being a PID means:

There exists a unique non-constant monic polynomial \( q_T \) in \( F[t] \) satisfying \( \ker \varphi = (q_T) > 0 \).

We call \( q_T \) the minimal polynomial of \( T \) over \( F \).

**Note:** By definition, \( q_T(T) = 0 \). (This polynomial need not be irreducible.)

The vector space \( V \) is an \( \text{End}_F(V) \)-module via evaluation so becomes an \( F[t] \)-module via the pullback along \( \varphi \), i.e., \( f \cdot v := \varphi(f)(v) = f(T)(v) \) for all \( f \in F[t] \) and all \( v \in V \). Thus \( V \) is an \( F[t] \)-module via evaluation at \( T \). As \( V \) is a finite dimensional vector space over \( F \), it is a finitely generated \( F[t] \)-module. Moreover, \( q_T \neq 0 \) in \( F[t] \) and \( q_T \cdot v = q_T(T)(v) = 0v = 0 \) for all \( v \) in \( V \), so \( V \) is a finitely generated torsion \( F[t] \)-module. Therefore, the two forms of the Fundamental Theorem of fg Modules over a PID are applicable. We shall use both forms. To use them, we need to investigate cyclic submodules of the \( F[t] \)-module \( V \).

We shall first apply the first form of the Fundamental Theorem 43.6 to the construction above. We shall use it for the case of an arbitrary field. We begin with some preliminaries.

**Lemma 44.2.** Suppose \( V \) is a finite dimensional vector space over \( F \) and \( T : V \rightarrow V \) a linear operator. Let \( V \) be an \( F[t] \)-module by evaluation at \( T \). Then

\[
(q_T) = \text{ann}_{F[t]} V := \{ f \in F[t] \mid f v = 0 \text{ for all } v \in V \} = \bigcap_{v \in V} \text{ann}_{F[t]} v.
\]

**Proof.** We know that \( q_T v \) is zero for all \( v \) in \( V \), so \( q_T \) lies in \( \text{ann}_{F[t]} V \). Let \( \varphi : F[t] \rightarrow \text{End}_F(V) \) be evaluation at \( T \). Suppose that \( f \in \text{ann}_{F[t]} V \). Then \( 0 = f v = f(T)(v) \) for all \( v \in V \), hence \( f(T) \) is the zero endomorphism in \( \text{End}_F(V) \), so \( f \in \ker \varphi = (q_T) \). The result follows. \( \square \)

We need another fact from linear algebra.

**Definition 44.3.** Let \( F \) be a field and \( h = t^d + a_{d-1}t^{d-1} + \cdots + a_1 t + a_0 \) be a monic polynomial in \( F[t] \) of positive degree \( d \). Define the
companion matrix of $h$ to be the matrix in $\mathbb{M}_d(F)$ given by
\[
C_h := \begin{pmatrix}
0 & 0 & \cdots & 0 & -a_0 \\
1 & 0 & \cdots & 0 & -a_1 \\
0 & 1 & \cdots & 0 & -a_2 \\
\vdots & \vdots & & \vdots & \vdots \\
0 & 0 & \cdots & 1 & -a_{d-1}
\end{pmatrix}.
\]

Remarks 44.4. Let $F$ be a field and $g, h$ be monic non-constant polynomials in $F[t]$.
1. $g = h$ if and only if $C_g = C_h$,
2. $h = f_{C_h}$, the characteristic polynomial of the matrix $C_h$:
   If $h = t^d + a_{d-1}t^{d-1} + \cdots + a_0$ with $d > 0$, then
   \[
f_{C_h} = \det \begin{pmatrix}
t & 0 & \cdots & 0 & a_0 \\
-1 & t & \cdots & 0 & a_1 \\
0 & -1 & \cdots & 0 & a_2 \\
\vdots & \vdots & & \vdots & \vdots \\
0 & 0 & \cdots & -1 & t + a_{d-1}
\end{pmatrix}.
\]

Using induction and expanding the determinant by minors along the top row yield
\[
f_{C_h} = t(t^{d-1} + a_{d-1}t^{d-2} + \cdots + a_1) + (-1)^{d-1}a_0(-1)^{d-1} = h.
\]

To apply the Fundamental Theorem 43.6 to a finite dimensional vector space over a field $F$ and linear operator $T$ with $F[t]$-module structure given by evaluation at $T$, we must compute $F[t]$-cyclic submodules. The key is to determine an appropriate $F$-basis, for such a submodule.

Proposition 44.5. Suppose $V$ is a finite dimensional vector space over $F$ and $T : V \to V$ a linear operator. Let $V$ be a cyclic $F[t]$-module by evaluation at $T$, i.e., there exists a vector $v$ in $V$ satisfying $V = F[t]v$. Let $B = \{v, T(v), \ldots, T^{d-1}(v)\}$ with $d$ the degree of the minimal polynomial $q_T$ of $T$. Then
1. $B$ is an ordered basis for $V$.
2. $\text{ann}_{F[t]} v = \text{ann}_{F[t]} V = (q_T)$
3. The matrix representation $[T]_B$ of $T$ in the $B$ basis satisfies $[T]_B = C_{q_T}$.
4. The minimal polynomial $q_T$ satisfies $q_T = f_T$.

Proof. $B$ spans $V$ as a vector space over $F$: Let $w$ be a vector in $V = F[t]v$, so there exists a polynomial $h$ in $F[t]$ satisfying $w = hv = h(T)(v)$, As $F[t]$ is a euclidean domain with euclidean
function the degree, we can write \( h = q_T s + r \) for some polynomials \( s \) and \( r \) in \( F[t] \) with \( r = 0 \) or \( \deg r < \deg q_T \). Evaluating this polynomial at \( T \) yields

\[
w = hv = (q_T s + r)v = q_T(T)s(T)(v) + r(T)(v) = r(T)(v)
\]

which lies in the span of \( B \).

**B is linearly independent (hence is a basis for \( V \)):** Suppose that \( \sum_{i=0}^{d-1} a_i T^i(v) = 0 \) for some \( a_i \) in \( F \). Set \( g = \sum_{i=0}^{d-1} a_i t^i \) in \( F[t] \). Then \( gw = g(T)(v) = 0 \), so \( gfv = 0 \) in \( V = F[t]v \) for all polynomials \( f \) in \( F[t] \), i.e., \( gw = 0 \) for all \( w \) in \( V \). By Lemma 44.2, we have \( g \in \text{ann}_{F[t]} V = \text{ann}_{F[t]} v = (q_T) \), so \( q_T | g \). As \( g = 0 \) or \( \deg g < \deg q_T \), we must have \( g = 0 \), i.e., \( a_i = 0 \) for all \( i \). The proof of (3) is simple and (4) follows from Remarks 44.4. \( \square \)

We now apply the Fundamental Theorem 43.6.

**Theorem 44.6.** (Rational Canonical Form) **Suppose** \( V \) **is a nontrivial finite dimensional vector space over \( F \) and** \( T : V \rightarrow V \) **a linear operator. Let** \( V \) **be an \( F[t] \)-module by evaluation at \( T \). Then there exists a direct sum decomposition of \( F[t] \)-modules**

\[
V = V_1 \oplus \cdots \oplus V_r
\]

for some \( r \) with each \( V_i \) a cyclic \( F[t] \)-module, \( 1 \leq i \leq r \), and satisfying:

1. Let \( q_i \) be the monic non-constant polynomial in \( F[t] \) satisfying \( (q_i) = \text{ann}_{F[t]} V_i \) for \( 1 \leq i \leq r \). Then the \( q_i \) satisfy \( q_1 | \cdots | q_r \) in \( F[t] \) and are unique relative to this ordering and are invariant factors of \( V \) as an \( F[t] \)-module.
2. The polynomial \( q_r \) in (1) is the minimal polynomial \( q_T \) of \( T \), and \( (q_T) = \text{ann}_{F[t]} V \).
3. Let \( q_i, i = 1, \ldots, r \) be as in (1). Then \( f_T = q_1 \cdots q_r \).
4. There exist ordered bases \( B_i \) for \( V_i, 1 \leq i \leq r \) satisfying:
   - (a) \( T|_{V_i} \) lies in \( \text{End}_F(V_i) \) for \( i = 1, \ldots, r \).
   - (b) \( [T|_{V_i}] g_i = C_{q_i} \) for \( i = 1, \ldots, r \).
   - (c) \( q_i = f_T|_{V_i} = q_T|_{V_i} \).
   - (d) Let \( B = B_1 \cup \cdots \cup B_r \) is an ordered basis for \( V \) and

\[
[T]_B = \begin{pmatrix}
C_{q_1} & \cdots & 0 \\
\vdots & \ddots & \vdots \\
0 & \cdots & C_{q_r}
\end{pmatrix} \quad \text{(in block form.)}
\]

This matrix is unique relative to the monic polynomials in (1) as ordered there.
The monic polynomials, $q_1, \ldots, q_r$ in (1) in the theorem are called the invariant factors of $T$ and the matrix in (4) is called the rational canonical form of $T$.

**Proof.** Apply the Fundamental Theorem 43.6 and Proposition 44.5 after noting that $V_i$ being $F[t]$-cyclic implies that $\text{ann}_{F[t]} V_i = (q_i) = (q_{|V_i})$ and $T|_{V_i} : V_i \to V_i$ as $V_i$ is an $F[t]$-module, i.e., $V_i$ is $T$-invariant as $tv = T(v)$ for all $v \in V_i$ for $i = 1, \ldots, r$. □

**Examples 44.7.** Suppose $V$ is a nontrivial finite dimensional vector space over $F$ and $T : V \to V$ a linear operator such that $T^m = 0$ for some $m > 0$ (we say that $T$ is nilpotent), then the rational canonical form of $T$ is upper triangular with diagonal entries all 0.

An immediate consequence of Theorem 44.6 is the following:

**Corollary 44.8.** Suppose $V$ is a nontrivial finite dimensional vector space over $F$ and $T : V \to V$ a linear operator. Let $V$ be an $F[t]$-module by evaluation at $T$. Then $q_T | f_T$ in $F[t]$. In particular, $q_T$ and $f_T$ have the same roots in $F$ and $q_T$ is a product of linear polynomials in $F[t]$ if and only if $f_T$ is a product of linear polynomials in $F[t]$.

Of course, the multiplicity of a root of $q_T$ and $f_T$ may be different. [Recall from linear algebra that the roots of $f_T$ are the eigenvalues of $T$.]

An important consequence of Theorem 44.6 is the well-known:

**Corollary 44.9.** (Cayley-Hamilton Theorem) Suppose $V$ is a nontrivial finite dimensional vector space over $F$ and $T : V \to V$ a linear operator. Then the characteristic polynomial, $f_T$ in $F[t]$, satisfies $f_T(T) = 0$, i.e., $f_T \in \text{ann}_{F[t]} V = (q_T)$. In particular, $q_T | f_T$ in $F[t]$ and $q_T$ and $f_T$ have the same irreducible factors (although not necessarily with the same multiplicities).

The rational canonical forms of a matrices in $\mathbb{M}_n(F)$ give an excellent system of representatives of similarity classes of such matrices. To establish this, we need the following observation.

**Observation 44.10.** If $R$ is a domain, $A$ a matrix in $\mathbb{M}_n(R)$ with $\det A$ nonzero, and $AX = 0$ ($X$ and 0 in $R^{n \times 1}$), then $X = 0$, i.e., $A$ is an $R$-monomorphism when viewed as a map $R^{n \times 1} \to R^{n \times 1}$. [This follows as $A$ is invertible in $\mathbb{M}_n(qf(R))$.]

**Theorem 44.11.** (Classification of Similarity Classes of Matrices over a Field) Let $F$ be a field and $A$ and $B$ be two $n \times n$ matrices in $\mathbb{M}_n(F)$.

(1) There exists a unique matrix $C$ in $\mathbb{M}_n(F)$ in rational canonical form with $A$ and $C$ similar.
(2) The following are equivalent:
   (a) $A \sim B$.
   (b) $A$ and $B$ have the same rational canonical form.
   (c) $A$ and $B$ have the same invariant factors.
   (d) $tI - A$ and $tI - B$ have the same Smith Normal Form when we choose monic invariants and the nonidentity monic invariant factors are the same as those for $A$ and $B$ in (c).
   (e) $tI - A$ and $tI - B$ are equivalent in $\mathbb{M}_n(F[t])$.

**Proof.** Let

$$S_{n,1} = \{e_1, \ldots, e_n\} \text{ with } e_i = (0 \ 0 \ 1 \ 0 \ \cdots \ 0)^t$$

be the standard basis for $F^{n \times 1}$ and also for $F[t]^{n \times 1}$ as free modules over $F$ and $F[t]^{n \times 1}$, respectively. [Aside: $C = \{t^j e_i \mid 1 \leq i \leq n, \ j \geq 0\}$ is an $F$-basis for $F[t]^{n \times 1}$. The linear operator $A : F^{n \times 1} \rightarrow F^{n \times 1}$ given by $v \mapsto Av$ satisfies $A = [A]_{S_{n,1}}$, so there exists a basis $B$ for $F^{n \times 1}$ such that $[A]_B$ is in rational canonical form by the Change of Basis Theorem (cf. Appendix B). Therefore, we have (1) and the equivalence of (a), (b), and (c) in (2) follow easily.

By the Universal Property of Free Modules 38.3, there exists a unique $F[t]$-endomorphism (hence also $F$-linear) $g : F[t]^{n \times 1} \rightarrow F[t]^{n \times 1}$ extending $e_j \mapsto te_j - Ae_j$. So $g = tI - A$. We, therefore, have a sequence

$$0 \rightarrow F[t]^{n \times 1} \xrightarrow{g} F[t]^{n \times 1} \xrightarrow{e_A} F^{n \times 1} \rightarrow 0 \quad (44.12)$$

with $e_A$ evaluation at $A$, i.e., defined by $fe_j \mapsto f(A)e_j$ for all $f \in F[t]$ and all $j$. We shall show that this sequence is an exact sequence.

As $\det(tI_n - A) = f_A$, the characteristic polynomial of the matrix $A$, hence nonzero, we have $g$ an $F[t]$-monomorphism by Observation 44.10 above. It is also clear that the evaluation map $e_A$ is a surjection. So we must show exactness at the middle term, i.e., $\text{im } g = \ker e_A$.

Clearly, $\text{im } g \subset \ker e_A$. Therefore, $x + \text{im } g \mapsto x + \ker e_A$ defines an an $F[t]$-epimorphism $F[t]^{n \times 1} / \text{im } g \rightarrow F[t]^{n \times 1} / \ker e_A$. As we want to show $\text{im } g = \ker e_A$, it suffices to show that this map is an isomorphism. Composing this map with $\overline{e_A} : F[t]^{n \times 1} / \ker e_A \rightarrow F^{n \times 1}$, the map induced by $e_A$ given by the First Isomorphism Theorem, defines an $F[t]$-epimorphism $\overline{e_A} : \text{coker } g \rightarrow F^{n \times 1}$, so it suffices to show that $\overline{e_A}$ is an isomorphism, i.e., a bijection.

By Remark 43.10(1), a Smith Normal Form $S$ of the matrix representation $tI_n - A$ of $g$ determines $\text{coker } g$ as an $F[t]$-module. We may assume that diagonal matrix $S$ was chosen with its nonzero diagonal
entries monic polynomials. As \( \det(tI_n - A) = f_A \neq 0 \), \( S \) has no nonzero diagonal entries.

Let the non-constant monic diagonal entries of \( S \) be \( d_1 | \cdots | d_r \). We know that \( S = P[g]_{S, n}Q \) with \( P, Q \in \GL_n(F[t]) \), and an invertible matrix in \( \GL_n(F[t]) \) has determinant in \( F[t]^\times = F^\times \). (This also follows from the algorithm given in Appendix C, Theorem C.2, since Type I, II, and III matrices have determinant in \( F^\times \)). Consequently, we have \( f_A = d_1 \cdots d_r \) as \( f_A \) and all the \( d_i \) are monic. Hence

\[
\dim_F \coker g = \sum_{i=1}^r \dim_F F[t]/(d_i)
\]

\[
= \sum_{i=1}^r \deg d_i = \deg f_A = n = \dim_F F^{n \times 1}
\]

using the isomorphism we showed when proving the Fundamental Theorem 43.6 and Proposition 33.14. It follows that the \( F[t] \)-epimorphism hence \( F \)-epimorphism \( e_A \) is an \( F \)-isomorphism, hence \( \im(g) = \ker e_A \) and the sequence in 44.12 is exact.

It follows from the exactness of 44.12 that the non-constant invariant factors of a Smith Normal Form \( S \) determines the invariants of \( F^{n \times 1} \) as an \( F[t] \)-module which therefore must be the same as those determined by the invariants of \( A \). The equivalence of (c), (d), and (e) follow.

\[\square\]

**Remark 44.13.** As pointed out in the proof above, the Smith Normal Form for a matrix \( A \in \M_n(F) \) is determined by the Smith Normal Form of the matrix \( tI - A \in \M_n(F[t]) \). In particular, the Smith Normal form of a linear operator on a finite dimensional \( F \)-vector space can be computed using the algorithm given in Appendix C, Theorem C.2, hence the invariant factors of \( A \).

**Corollary 44.14.** Let \( F \) be a field and \( A \in \M_n(F) \). Then \( A \) is similar to its transpose \( A^t \).

**Proof.** It follows by Theorem 42.8 (or by the algorithm given in the proof of Theorem C.2 in Appendix C) that a matrix and its transpose in \( \M_n(F[t]) \) have equivalent Smith Normal Forms. In particular, the matrices \( tI - A \) and \( tI - A^t \) have the same Smith Normal Form when we choose monic invariants.

\[\square\]

Using the proof above, we can give a summary of what is going on. Instead of giving this summary using matrices, we use linear operators.
Summary 44.15. Let $V$ be a finite dimensional vector space over $F$ of dimension $n$ and $T : V \to V$ a linear operator. We view $V$ as an $F[t]$-module by $tv := T(v)$ for all $v \in V$. We then have an exact sequence
\begin{equation}
0 \to F[t]^n \xrightarrow{t1_F[t]^n - T} F[t]^n \xrightarrow{e_T} V \to 0
\end{equation}
with the $F[t]$-homomorphism $e_T$ defined by $\sum f_i e_i \mapsto \sum f(T)e_i$, where $e_i$ in the $i$th basis element in the (ordered) standard basis $S_n$ for $F[t]^n$ (and $F^n$). We call the sequence 44.16 the characteristic sequence of $T$. The matrix $t1_F[t]^n - [T]_{S_n}$ is called the characteristic matrix of $[T]_{S_n}$. It has nonzero determinant $f_T$, the characteristic polynomial of $T$ (and is the reason that $t1_F[t]^n - [T]_{S_n}$ is a monomorphism). In particular, the characteristic matrix of $[T]_{S_n}$ is nonsingular, so has a unique Smith Normal Form of $(tI^n - [T]_{S_n})$ given by diag($1$, $\ldots$, $q_1$, $\ldots$, $q_r$). (The entries of a Smith Normal Form are unique up to units.)

As the sequence 44.16 is exact, we have
\[
coker(t1_F[t]^n - T) = F[t]^n / \text{im}(t1_F[t]^n - T) = F[t]^n / \ker(t1_F[t]^n - T) \cong V.
\]
By the Fundamental Theorem, we can decompose the finitely generated torsion $F[t]$-module $V$ into cyclic $F[t]$-submodules, say
\[
V = V_1 \coprod \cdots \coprod V_r
\]
with $V_i = F[t]v_i$ some $v_i$ for each $i$ and satisfying
\[
V_i \cong F[t] / \text{ann}_{F[t]}(v_i) = F[t] / \text{ann}_{F[t]}(V_i) = F[t] / (q_i), \quad i = 1, \ldots, r
\]
with the descending chain of ideals
\[
(q_1) \supset \cdots \supset (q_r)
\]
unique. Moreover,
\[
q_i = q_T|v_i; \quad q_r = q_T; \quad f_T = q_1 \cdots q_r,
\]
and
\[
RCF(T) := RCF([T]_{S_n}) = \begin{pmatrix}
C_{q_1} & & \\
& \ddots & \\
& & C_{q_r}
\end{pmatrix}.
\]
In particular, to find the invariant factors of $T$ one computes the Smith Normal Form of the characteristic matrix using the algorithm given in the Appendix C.

We turn to the second form of the Fundamental Theorem. It is immediate that we have the following:
Theorem 44.17. Suppose $V$ is a nontrivial finite dimensional vector space over $F$ and $T : V \rightarrow V$ a linear operator. Let $V$ be an $F[t]$-module by evaluation at $T$. Then there exists a direct sum decomposition of $F[t]$-modules $V = V_{11} \oplus \cdots \oplus V_{1r_1} \oplus \cdots \oplus V_{r_1} \oplus \cdots \oplus V_{rn_r}$ for some $r$ and $r_i$ with each $V_{ij}$ a cyclic $F[t]$-module, $1 \leq i \leq r$ and $1 \leq j \leq i$, satisfying:

1. $V_{ij} \cong F[t]/(q_{ij})$ with $q_{ij} \in F[t]$ monic.
2. $1 \leq e_{i1} \leq e_{i2} \leq \cdots \leq e_{in_i}$.
3. $(q_{ij}) = \text{ann}_{F[t]} V_{ij}$ for $1 \leq i \leq r$ and $1 \leq e_{i1} \leq e_{i2} \leq \cdots \leq e_{in_i}$, are unique relative to the ordering in (2), and are the elementary divisors of $V$ as an $F[t]$-module.
4. The polynomial $q_{ij}$ satisfy $q_{ij} = q_T|_{V_{ij}} = f_T|_{V_{ij}}$ for $1 \leq i \leq r$.
5. $f_T = \prod_{i,j} q_{ij}$.
6. There exist ordered bases $B_{ij}$ for $V_{ij}$, $1 \leq i \leq r$ and $1 \leq e_{i1} \leq e_{i2} \leq \cdots \leq e_{in_i}$ satisfying:
   a. $T|_{V_{ij}}$ lies in $\text{End}_F(V_{ij})$ for $i = 1, \ldots, r$ and $1 \leq e_{i1} \leq e_{i2} \leq \cdots \leq e_{in_i}$.
   b. $[T|_{V_{ij}}]_{B_{ij}} = C_{q_{ij}}$, for $i = 1, \ldots, r$ and $1 \leq e_{i1} \leq e_{i2} \leq \cdots \leq e_{in_i}$.
   c. Let $B = B_1 \cup \cdots \cup B_r$. Then $B$ is an ordered basis for $V$, and

$$[T]_B = \begin{pmatrix} C_{q_{11}} & \cdots & 0 \\ \vdots & \ddots & \vdots \\ 0 & \cdots & C_{q_{rn_r}} \end{pmatrix} \quad \text{(in block form.)}$$

This matrix is unique relative to the monic polynomials in (1) as ordered in (2) up to a permutation of the blocks determined by $\prod_{i=1}^{n_i} q_{ij}$.

We now shall apply the second form for matrices in the more interesting case when the field is an algebraically closed field, i.e., when non-constant polynomials always factor into a product of linear polynomials, as it is this case that the elementary divisors arising from this form are useful. Instead of the form above, we shall get a more useful form. We first establish the cyclic case when the minimal polynomial is a power of a linear polynomial.

Lemma 44.18. Suppose $V$ is a finite dimensional vector space over $F$ and $T : V \rightarrow V$ a linear operator. Let $V$ be a cyclic $F[t]$-module by evaluation at $T$, i.e., there exists a vector $v$ in $V$ satisfying $V = F[t]v$. 
Suppose that the minimal polynomial $q_T$ of $T$ is $(t-a)^d$ in $F[t]$. Then $B = \{v, (T-a)(v), \ldots, (T-a)^{d-1}(v)\}$ is an ordered basis for $V$ and

$$[T]_B = \begin{pmatrix} a & 0 & 0 & \cdots & 0 \\ 1 & a & 0 & \cdots \\ 0 & 1 & a & \cdots \\ \vdots & \ddots & \ddots & \ddots \\ 0 & \cdots & 1 & a \end{pmatrix} \text{ in } M_d(F).$$

**Proof.** We know that $\mathcal{C} = \{v, T(v), \ldots, T^{d-1}(v)\}$ is an $F$-basis for $V$ by our previous work. If $B$ is linearly dependent, then there exists an equation

\[0 = \sum_{i=0}^{d-1} a_i (T-a)^i v \text{ for some } a_i \text{ in } F, \text{ not all zero.}\]  

Choose $N$ maximal such that $a_N \neq 0$ and expand (*) to get

\[0 = a_N T^N v + \sum_{i=0}^{N-1} b_i T^i v \text{ for some } b_i \text{ in } F.\]

As $\mathcal{C}$ is linearly independent, $a_N = 0$, a contradiction. Therefore, $B$ is linearly independent hence a basis as $|B| = |\mathcal{C}|$. We have

\[0 = q_r(T) = (T-a1_V)^r \text{ and} \]

\[T(T-a1_V)^i = ((T-a1_V) + a)(T-a1_V)^i = (T-a1_V)^{i+1} + a(T-a1_V)^i \]

for $i = 0, \ldots, r-1$. So

\[T(T-a)^i(v) = (T-a)^{i+1}(v) + a(T-a)^i(v) \text{ for } i = 0, \ldots, r-1\]

and the last statement follows.

When the minimal polynomial of a linear operator on a finite dimensional vector space factors into a product of linear terms, the Fundamental Theorem 43.22 and the last lemma now immediately implies the following:

**Theorem 44.19.** (Jordan Canonical Form) Suppose $V$ is a nontrivial finite dimensional vector space over $F$ and $T : V \rightarrow V$ a linear operator. Let $V$ be an $F[t]$-module by evaluation at $T$. Suppose that $q_T$ factors into linear terms in $F[t]$. Then there exists an ordered basis $B$
for $V$ satisfying

\[ [T]_B = \begin{pmatrix} A_1 & \cdots & 0 \\ \vdots & \ddots & \vdots \\ 0 & \cdots & A_r \end{pmatrix} \quad \text{with} \quad A_i = \begin{pmatrix} a_i & 0 & \cdots & 0 \\ 1 & a_i & 0 & \cdots \\ 0 & 1 & a_i & 0 \\ \vdots & \ddots & \ddots & \ddots \\ 0 & \cdots & 1 & a_i \end{pmatrix} \]

for $i = 1, \ldots, r$. The $a_i$, $i = 1, \ldots, r$ are the eigenvalues of $T$ (with $a_i = a_j$, $i \neq j$ possible). This matrix of these blocks is unique up to the order of the blocks. The monic polynomials $f_{A_i}$, $1 \leq i \leq r$, are the elementary divisors of $V$ with the above $F[t]$-module structure and the characteristic polynomial of $T$ satisfies $f_t = \prod_{i=1}^r f_{A_i}$.

In the theorem, the blocks $A_i$ are called Jordan blocks and, as the monic elementary divisors $f_{A_i}$ are unique, they are called the elementary divisors of $T$. The matrix representing $T$ in \ref{44.19} is called the Jordan canonical form of $T$.

**Corollary 44.20.** Let $A$ and $B$ lie in $\mathbb{M}_n(F)$. If each of the characteristic polynomials $f_A$ and $f_B$ factor into a product of linear terms in $F[t]$, then $A$ and $B$ are similar in $\mathbb{M}_n(F)$ if and only if $A$ and $B$ have the same Jordan canonical form (up to block order) if and only if they have the same elementary divisors (counted with multiplicity). In particular, if $F$ is an algebraically closed field, then a matrix in $\mathbb{M}_n(F)$, $n \in \mathbb{Z}^+$, is determined up to similarity by its Jordan Canonical Form.

**Remark 44.21.** The rational canonical form of a matrix is quite computable, as the division algorithm is computable, but the Jordan canonical form (when it exists) maybe harder to compute as it depends on factoring the invariant factors of the monic polynomials arising in the rational canonical form.

**Corollary 44.22.** Let $F$ be a field and $A$ a matrix in $\mathbb{M}_n(F)$. Then $A$ is similar to a diagonal matrix if and only if the minimal polynomial $q_A$ factors as a product of monic linear factors without repetition, i.e., $q_A = (t - \lambda_1) \cdots (t - \lambda_r)$ for some distinct $\lambda_1, \ldots, \lambda_r$ in $F$.

We say that a non-constant polynomial in $F[t]$ splits over $F$ if it is a product of linear polynomials in $F[t]$. The corollary says that the matrix is diagonalizable if and only if $q_A$ splits over $F$ without any multiple roots.
Corollary 44.23. Let $F$ be a field and $A$ a matrix in $\mathbb{M}_n(F)$. Then $A$ is similar to a triangular matrix if and only if the minimal polynomial $q_A$ splits over $F$.

We say that $A$ is triangularizable if $q_A$ splits over $F$.

Examples 44.24. Let $F$ be a field

1. Similarity classes of matrices in $\mathbb{M}_3(F)$ are completely determined by the minimal and characteristic polynomials of a matrix: Let $A$ lie in $\mathbb{M}_3(F)$ and $q_1 | \cdots | q_r$ be the invariant factors of $A$. So $f_A = q_1 \cdots q_r$ with $q_r = q_A$. As $A \in \mathbb{M}_3(F)$, we know that $r \leq 3$.

$r = 1$: We must have $q_1 = q_A = f_A$ determines $A$ up to similarity, so $A$ is similar to $C_{q_1} = C_{q_A}$. We can say more if $q_A$ is reducible. First suppose that $q_A = (t - \lambda)h$ in $F[t]$ with $h$ irreducible. By Exercise 44.25(2) below, $A$ is similar to the matrix

\[
\begin{pmatrix}
\lambda & 0 & 0 \\
0 & C_h \\
0 & 0 
\end{pmatrix},
\]

which is unique up to block order. If $q_A$ splits over $F$, say $q_A = (t - \lambda_1)(t - \lambda_2)(t - \lambda_3)$, then $A$ is similar to

\[
\begin{pmatrix}
\lambda_1 & 0 & 0 \\
0 & \lambda_2 & 0 \\
0 & 0 & \lambda_3 
\end{pmatrix}
\]

if $\lambda_1, \lambda_2, \lambda_3$ are distinct,

\[
\begin{pmatrix}
\lambda_1 & 0 & 0 \\
1 & \lambda_1 & 0 \\
0 & 0 & \lambda_3 
\end{pmatrix}
\]

if $\lambda_1 = \lambda_2 \neq \lambda_3$,

\[
\begin{pmatrix}
\lambda_1 & 0 & 0 \\
1 & \lambda_1 & 0 \\
0 & 1 & \lambda_1 
\end{pmatrix}
\]

matrices in Jordan Canonical Form, unique up to order of the blocks.

$r = 2$: We have $q_1 | q_2 = q_A$ and $f_A = q_1 q_2$, so $q_1 = f_A / q_A$ in $F[t]$. Therefore, $q_A$ and $f_A$ determine $A$ up to similarity, and $A$ is similar to

\[
\begin{pmatrix}
C_{q_1} & 0 \\
0 & C_{q_2} 
\end{pmatrix}.
\]

Note that if this is the case, $\deg q_1 = 1$ and $\deg q_2 = 2$, so $q_1 | q_2$ implies that $q_2$ splits so $A$ is triangularizable, i.e., a Jordan canonical
form for $A$ exists. If $q_1 = t - \lambda_1$ and $q_2 = (t - \lambda_1)(t - \lambda_2)$, then $A$ is similar to

$$
\begin{pmatrix}
\lambda_1 & 0 & 0 \\
0 & \lambda_1 & 0 \\
0 & 0 & \lambda_2
\end{pmatrix}
$$

if $\lambda_1 \neq \lambda_2$,

$$
\begin{pmatrix}
\lambda_1 & 0 & 0 \\
0 & \lambda_1 & 0 \\
0 & 1 & \lambda_1
\end{pmatrix}
$$

if $\lambda_1 = \lambda_2$,

unique up to the order of the blocks.

$r = 3$: We have $q_1 \mid q_2 \mid q_3 = q_A$, so $q_1 = q_2 = q_3 = q_A$. So in this case $q_A$ and $f_A$ determine $A$ up to similarity, and $A$ is diagonalizable with a single eigenvalue, i.e., $A$ is similar to the matrix $\lambda I$ with $\lambda$ an (the) eigenvalue of $A$. In particular, $A$ is similar to a matrix in Jordan canonical form.

2. Similarity classes of matrices in $M_4(F)$ are not determined by minimal and characteristic polynomials. For example the matrices

$$
A = 
\begin{pmatrix}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 1 & 1
\end{pmatrix}
$$

and

$$
B = 
\begin{pmatrix}
1 & 0 & 0 & 0 \\
1 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 1 & 1
\end{pmatrix}
$$

have

$q_A = (t - 1)^2$

$f_A = (t - 1)^4$

elementary divisors

t - 1, t - 1, (t - 1)^2$

invariant factors

t - 1 | t - 1 | (t - 1)^2

so they are not similar.

3. A system of representatives for the similarity classes of matrices in $M_4(F)$ having characteristic polynomial $t(t - 1)^3$ are:
4. A system of representatives for the similarity classes of matrices in $\mathbb{M}_4(F)$ with minimal polynomial $t(t-1)$ are:

- $t(t-1)$
- $t, (t-1), (t-1)^2$
- $t-1 | t-1 | t(t-1)^2$

as these matrices must be diagonalizable with eigenvalues 0 and 1.

5. Let $A$ be a $3 \times 3$ matrix in $\mathbb{M}_3(\mathbb{Q})$ satisfying $A^3 = I$. Then $q_A | t^3 - 1 = (t-1)(t^2 + t + 1)$ in $\mathbb{Q}[t]$, so the only possibilities are $t-1$, $t^2 + t + 1$, or $t^3 - 1$. If $q_A = t^2 + t + 1$, then $q_1 | q_2 = q_A$, but $q_A$ is irreducible, so no linear polynomial divides it, so this case cannot occur. If $q_A = t-1$, $A$ is diagonalizable and $t-1 | t-1 | t-1$ are the invariants. So the possible rational canonical forms for $A$ are:

- $\begin{pmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1 \end{pmatrix}$
- $\begin{pmatrix} 0 & 0 & 1 \\ 1 & 0 & 0 \\ 0 & 1 & 0 \end{pmatrix}$

[Note: Using Exercise 44.25(2) below, an alternate matrix similar to

$\begin{pmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1 \end{pmatrix}$

is $\begin{pmatrix} 1 & 0 & 0 \\ 0 & 0 & -1 \\ 0 & 1 & -1 \end{pmatrix}$]

Exercises 44.25.

1. Prove Observation 44.10.

2. Let $F$ be a field and $A$ a matrix in $\mathbb{M}_n(F)$. Let $f_1, \ldots, f_r$ be the monic elementary divisors of $A$, i.e., those of $F^{n \times 1}$ as a $F[t]$ module via $tv = Av$ for all $v$ in $F^{n \times 1}$. Show that $A$ is similar to the matrix
3. Let $A$ be a matrix in $\mathbb{M}_5(\mathbb{Q})$ satisfying $A^3 = I$. Describe up to similarity all possible rational canonical forms of $A$ and justify your answer.

4. Describe, up to similarity, all $4 \times 4$ matrices $A$ over a field $F$ such that $A^5 = I$, in the following three cases and justify your answer.
   (i) $F$ is the field of rational numbers.
   (ii) $F$ is $\mathbb{Z}/2\mathbb{Z}$
   (iii) $F$ is $\mathbb{Z}/5\mathbb{Z}$

5. Compute the Jordan and Rational Canonical Forms of the matrix

\[
\begin{pmatrix}
0 & 0 & 0 & 0 \\
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
1 & 0 & 0 & 0
\end{pmatrix}
\]

6. Let $V$ be a real vector space, $T$ an $\mathbb{R}$-endomorphism. Suppose that the minimal polynomial of $T$ factors into linear terms over $\mathbb{R}$ with no repeated and no negative roots. Show that there exists an $\mathbb{R}$-endomorphism $S$ of $V$ such that $S^2 = T$.

7. Recall that a matrix $C \in \mathbb{M}_n(\mathbb{C})$ is called diagonalizable if there exists a matrix $Q \in \text{GL}_n(\mathbb{C})$ such that $QAQ^{-1}$ is a diagonal matrix. Show if $A, B \in \mathbb{M}_n(\mathbb{C})$ are both diagonalizable, then there is a matrix $P \in \text{GL}_n(\mathbb{C})$ satisfying both $PAP^{-1}$ and $PBP^{-1}$ are diagonal matrices if and only if $AB = BA$.

8. Determine a theorem for a canonical form for matrices in $\mathbb{M}_n(\mathbb{F})$ when $\mathbb{F}$ is an arbitrary field given by the Fundamental Theorem 43.22.

45. Addendum: Jordan Decomposition

We prove a weaker form of Jordan canonical form that only needs the Primary Decomposition Theorem, not the full Fundamental Theorem 43.6. [The Cayley-Hamilton Theorem can be proven without the full version of the Fundamental Theorem.]

**Definition 45.1.** Suppose that $\mathbb{F}$ is a field, $V$ a vector space over $\mathbb{F}$, and $T$ is a linear operator on $V$. Let

\[
E_T(\lambda) := \{v \in V \mid T(v) = \lambda v\},
\]
a subspace of $V$. If $\lambda$ is an eigenvalue of $T$, then $E_T(\lambda)$ is called the eigenspace of $\lambda$ relative to $T$, i.e., $E_T(\lambda)$ is an eigenspace if and only if it is not the zero subspace. We say that $T$ is semisimple if $V$ has a basis consisting of eigenvectors of $T$, equivalently, $V = \bigoplus_F E_T(\lambda)$. Of course, $T|_{E_T(\lambda)} = \lambda 1_{E_T(\lambda)}$ for all $\lambda \in F$.

**Observation 45.2.** Let $V$ be a vector space over $F$ and $S, T$ two commuting linear operators on $V$, i.e., $T \circ S = S \circ T$. If $\lambda$ is an eigenvalue of $T$, then $E_T(\lambda)$ is $S$-invariant, i.e., $S|_{E_T(\lambda)} : E_T(\lambda) \to E_T(\lambda)$:

If $v \in E_T(\lambda)$, then $T(S(v)) = S(T(v)) = S(\lambda v) = \lambda S(v)$.

The observation has the following consequence:

**Theorem 45.3.** Let $V$ be a finite dimensional vector space over the field $F$ and $S$ a set of commuting semisimple linear operators in $\text{End}_F(V)$, i.e., if $T$ and $S$ lie in $S$, they are semisimple and $ST = TS$. Then there exists a basis $B$ for $V$ consisting of eigenvectors for every $T$ in $S$ simultaneously. In particular, $T + S$ is semisimple for all $T, S \in S$.

**Proof.** If $T = \rho 1_V$ for some $\rho \in F$ for every $T \in S$, the result is trivial. So we may assume that there exists a $T$ in $S$ so that this is not true. As $T$ can have only finitely many eigenvalues, say $\lambda_1, \ldots, \lambda_r$, we have $V = E_T(\lambda_1) \oplus \cdots \oplus E_T(\lambda_r)$ with each $E_T(\lambda_j) < V$. By the observation, $S|_{E_T(\lambda_j)} \in \text{End}_F(E_T(\lambda_j))$ for all $S \in S$. As $q_S|_{E_T(\lambda_j)} | q_S$ (why?), we know that $q_S|_{E_T(\lambda_j)}$ splits over $F$, so $S|_{E_T(\lambda_j)}$ is semisimple.

By induction, the result holds on each $E_T(\lambda_j)$, hence on $V$. \hfill $\Box$

**Definition 45.4.** Let $F$ be a field, $V$ a nontrivial vector space over $F$, and $T$ a linear operator on $V$. We say that $T$ is nilpotent if there exists a positive integer $N$ such that $T^N = 0$ and unipotent if $T - 1_V$ is nilpotent.

Note that a unipotent linear operator is always an isomorphism (why?) and an operator that is both nilpotent and semisimple must be the zero linear operator.

**Theorem 45.5.** Let $V$ a finite dimensional vector space over $F$, and $T$ a linear operator on $V$. Suppose that $f_T$ splits over $F$. Then there exist unique linear operators $T_s$ and $T_n$ on $V$ semisimple and nilpotent respectively and if $T$ is a linear automorphism, there exists a unique unipotent operator $T_u$ satisfying the following:

1. (Additive Jordan Decomposition) $T = T_s + T_n$ and $T_sT_n = T_nT_s$.
2. (Multiplicative Jordan Decomposition) If $T$ is a linear automorphism, then $T = T_sT_u$ and $T_sT_u = T_uT_s$. 

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(3) There exist polynomials $g$ and $h$ in $F[t]$ with zero constant term, i.e., $g(0) = 0 = h(0)$ satisfying $T_s = g(T)$ and $T_n = h(T)$. In particular, if a linear operator $S$ on $V$ commutes with $T$, then it commutes with both $T_s$ and $T_n$.

(4) If $T$ is a linear automorphism and $T$ commutes with a linear operator $S$ on $V$, then it commutes with $T_s$ and $T_n$.

PROOF. Let $V$ be an $F[t]$-module by evaluation at $T$ and suppose that the characteristic polynomial of $T$ is $f_T = \prod_{i=1}^{r}(t - \lambda_i)^{e_i}$ in $F[t]$ with $\lambda_i \neq \lambda_j$ if $i \neq j$, and $e_i > 0$ for all $i$. Then by the Cayley-Hamilton Theorem, $f_T v = 0$ for all $v$ in $V$. Let

$$V_i = \ker(T - \lambda_i 1_V)^{e_i} \cong F[t]^{e_i}/((t - \lambda_i)^{e_i})$$

for $i = 1, \ldots, r$, so $V = V_1 \oplus \cdots \oplus V_r$ as $F[t]$-modules by the Primary Decomposition Theorem 43.20 (Why?). By the Chinese Remainder Theorem (the proof for $\mathbb{Z}$ works for any Euclidean domain — or follows from the general version as the ideal $(t - \lambda_i)^{e_i} + \sum_{j \neq i} (t - \lambda_j)^{e_j}$ is $F[t]$, there exists a polynomial $g$ in $F[t]$ satisfying $g \equiv \lambda_i \mod (t - \lambda_i)^{e_i}$, $1 \leq i \leq r$ and $g \equiv 0 \mod t$. (Of course, if some $\lambda_i = 0$, the last congruence is redundant.) In particular, $g(0) = 0$. Let $T_s = g(T)$.

Each $V_i$ is $T$-invariant, as it is an $F[t]$-module by evaluation at $T$, so each $V_i$ is also $T_s$-invariant. We also have $T_s|_{V_i} = \lambda_i 1_{V_i}$ for $i = 1, \ldots, r$, as the $\lambda_i$'s are eigenvalues of $T$. Therefore, $T_s$ is semisimple.

Let $h = t - g$ in $F[t]$, so $h(0) = 0$. Set $T_n = h(T)$. We have

$$T = g(T) + h(T) = T_s + T_n \quad \text{and} \quad T_s T_n = T_n T_s.$$

In addition,

$$0 = (T|_{V_i} - \lambda_i 1_{V_i})^{e_i} = (T|_{V_i} - T_n|_{V_i})^{e_i} = T_n^{e_i}|_{V_i}$$

for each $i$, so $T_n$ is nilpotent. This proves (3) and establishes the existence in (1).

If $T$ is a linear automorphism, then $0$ is not an eigenvalue of $T$, so $T_s$ is invertible. Let $T_u = 1_V + T_s^{-1}T_n$. Then $T_u$ is unipotent as $T_s^{-1}T_n = T_n T_s^{-1}$ is nilpotent. Statement (4) and the existence in (2) now follow.

We next show the uniqueness in (1). Suppose that we have a semisimple operator $T_s'$ and a nilpotent operator $T_n'$ satisfying

$$T_s + T_n = T_s' + T_n' \quad \text{and} \quad T_s' T_n' = T_n' T_s'$$

Then $T_s - T_s' = T_n' - T_n$. By (3), $T_s'$ and $T_n'$ commute with both $T_s$ and $T_n$, as they commute with $T = T_s' + T_n'$. Therefore, $T_n' - T_n$ is nilpotent and by Theorem 45.3 we have $T_n' - T_n'$ is semisimple. Therefore, $T_s - T_s' = T_n - T_n = 0$, i.e., $T_s = T_s'$ and $T_n = T_n'$. 

Finally, we show the uniqueness in (3). Suppose that
\[ T_s T_u = T'_s T'_u \quad \text{and} \quad T'_s T'_u = T'_n T'_u. \]
with \( T_u \) unipotent. Then \( S = T'_u - 1_V \) is nilpotent and commutes with \( T'_s \), hence \( T'_s S \) is nilpotent and \( T = T'_s + T'_s S \) is an additive Jordan decomposition, so \( T'_s = T_s \) (and \( T_n = T'_s S \)). Hence \( T_u = T^{-1}_s T = T'_s - T_u \), and we are done. \( \square \)

Our definition of a semisimple operator \( T \) on \( V \) is equivalent to the minimal polynomial \( q_T \) splitting without any repeated factors over \( F \). The theorem generalizes if we can find an alternate condition of a linear operator on \( V \) when \( q_T \) does not split over \( F \) that will agree with our definition of a semisimple operator over a field containing \( F \) in which \( q_T \) splits. Then we can mimic the proof on a factorization of \( q_T \), if it is not a product of distinct linear terms. To do this, one can redefine an operator \( T \) to be a semisimple operator on a vector space \( V \) over \( F \) if it satisfies the following: Let \( V \) be an \( F[t] \)-module by evaluation at \( T \). If \( W \) is an \( F[t] \)-submodule of \( V \) then \( V = W \oplus W' \) for some \( F[t] \)-module \( W' \), i.e., and exact sequence
\[ 0 \to W \to V \to W' \to 0 \]
of \( F[t] \)-modules is split exact (cf. Exercise 38.12(12)). With this definition, the theorem still holds if the characteristic of \( F \) is zero. [If the characteristic of \( F \) is positive, irreducible polynomials can have multiple roots over some field containing \( F \).]

**Exercises 45.6.**

1. Show a unipotent operator on a vector space is an isomorphism and an operator that is both nilpotent and semisimple on a vector space is the trivial operator.

2. In the proof of Theorem 45.5, show that \( V_i = \ker(T - \lambda_i 1_V)^{e_i} \cong F[t]/((t - \lambda_i)^{e_i}) \) is the \( (t - \lambda_i) \)-primary submodule of \( V \) as an \( F[t] \)-module.

3. In Theorem 45.5, show if \( W \) is an \( F[t] \)-submodule of \( V \), i.e., is \( T \)-invariant, then the Jordan decomposition(s) of \( T \) induce those on \( T|_W \) and \( T_{V/W} \) (where \( T_{V/W}(x + W) = T(x) \)).

46. Addendum: Cayley-Hamilton Theorem

In this addendum, we generalize the Cayley-Hamilton Theorem. Let \( R \) be a commutative ring and \( M \) a finitely generated \( R \)-module.
Let 

\[ M[t] := \{ \sum t^ix_i \mid x_i \in M \text{ almost all } x_i = 0 \} = \prod_{i=0}^{\infty} t^iM. \]

This is clearly an \( R \)-module and becomes a finitely generated \( R[t] \)-module in the obvious way.

For example, if \( V \) is a finite dimensional vector space over \( F \) on basis \( \{ v_i \} \), then \( V[t] \) is the \( F \)-vector space on basis \( \{ t^iv_j \mid i \geq 0, j \} \) and a free \( F[t] \)-module on basis \( \{ v_i \} \). So if \( \dim_F V = n, V[t] \cong F[t]^n \).

If \( f : M \to N \) is an \( R \)-homomorphism, define \( f[t] : M[t] \to N[t] \) by \( \sum t^i x_i \mapsto \sum t^i f(x_i) \), an \( R[t] \)-homomorphism.

Suppose that \( N \) is an \( R[t] \)-module. Then it is also an \( R \)-module and \( N \) is completely determined by

(i) the \( R \)-module structure.
(ii) The \( R \)-endomorphism \( f : N \to N \) given by \( f(x) = tx \).

In particular, an \( R \)-module becomes an \( R[t] \)-module via (i) and (ii).

We denote this \( R[t] \)-module \( N_f \). We then have an \( R[t] \)-epimorphism \( \varphi_f : N[t] \to N_f \) given by \( \sum t^i x_i \mapsto \sum f^i(x_i) \).

**Theorem 46.1.** Let \( R \) be a commutative ring and \( M \) a finitely generated \( R \)-module. Suppose that \( f \) is an \( R \)-endomorphism of \( M \). Then, in the notation above, we have an exact sequence

\[ 0 \to M[t] \xrightarrow{t1_{M[t]} - f[t]} M[t] \xrightarrow{\varphi_f} M_f \to 0 \]

of \( R[t] \)-modules.

The exact sequence in the theorem is called the characteristic sequence for \( f \).

**Proof.** Exactness at the middle \( M[t] \): As

\[
(\varphi_f \circ (t1_{M[t]} - f[t]))(\sum t^i x_i) = \varphi_f(\sum t^{i+1} x_i - t^i f(x_i)) \\
= \sum f^{i+1}(x_i) - f^{i+1}(x_i) = 0,
\]

we have \( \text{im}(t1_{M[t]} - f[t]) \subset \ker \varphi_f \). If \( x = \sum t^i x_i \in \ker f \), then \( \sum f^i(x_i) = 0 \), hence

\[ x = x - \sum f^i(x_i) = \sum_{i \geq 0} (t^i x_i - f^i(x_i)) = \sum_{i > 0} (t^i1_{M[t]} - f^i)(x_i). \]

Check if \( h_i = \sum_{j=0}^{i} t^{j}1_{M[t]}f^{i-j}[t] \), then

\[ x = \sum_{i > 0} (t1_{M[t]} - f[t])(\sum_{i > 0} h_i(x_i)) \text{ lies in } \text{im}(f(t1_{M[t]} - f[t])). \]
\( \varphi_f \) is surjective: This follows as \( \varphi_f(x) = x \) for all \( x \in M \).

\( t1_{M[t]} - f[t] \) is monic: We leave this part of the proof as an exercise noting that \( t1_{M[t]} - f[t] \) raises the “degree” by one and preserves the “leading coefficients”, so induction on degree works. \( \square \)

Cramer’s Rule can be stated in the following way: Let \( R \) be a commutative ring, \( f \) an \( R \)-endomorphism of the free \( R \)-module \( R^n \). Set \( M = \text{coker } f = R^n / \text{im } f \). Then there exists an \( R \)-endomorphism \( g \) of \( R^n \) such that \( fg = gf = (\det f) 1_{R^n} \). In particular, \( R^n \det f \subset \text{im } f \).

If \( M \) is a free \( R \)-module of rank \( n \) on basis \( \mathcal{B} \) and \( f \in \text{End}_R(M) \), then \( M[t] \) is \( R[t] \)-free on \( \mathcal{B} \) and \( [f]_\mathcal{B} = [f[t]]_\mathcal{B} \). If follows that \( P(t) := \det(t1_{M[t]} - f[t]) \) is that characteristic polynomial of \( f \). Applying the characteristic sequence of \( f \) in the theorem implies that \( P(f)M_f = 0 \). As the \( R \)-endomorphism of \( M_f \) defined by \( P_f(t) \) is just \( P_f(f) \), we have \( P_f(f) = 0 \), i.e., the Cayley-Hamilton Theorem.

**Exercise 46.2.** Prove that \( t1_{M[t]} - f[t] \) in the characteristic sequence is monic.
Part 5

Field Theory
CHAPTER XI

Field Extensions

In this chapter, we begin our study of fields. Given a non-constant polynomial, it is natural to ask what its roots are. To make this more precise, if $R$ is a ring and $f$ is a non-constant polynomial in a single variable over $R$, one asks does $f$ have any roots in $R$. If so how many and if not does there exist a ring $S$ with $R$ a subring such that it does. We have seen if $R$ is a domain, then the number of distinct roots $f$ is bounded by the degree of $f$. The first domain one encounters is the integers. The arithmetic problem, i.e., for $f$ in $\mathbb{Z}[t]$ and monic leads to the study of algebraic number theory that we shall investigate in Chapter XV and is harder than that over fields. If $R = F$ is a field, we already know how to obtain a field $K$ containing $F$ such that $f$ has a root in $K$. By induction, it is easy to construct a field $E$ containing $F$ for which the polynomial $f$ splits. We shall show the smallest such $E$ is essentially unique. Generalizing this construction, we indicate, using Zorn’s Lemma, how to do this for sets of polynomials in $F[t]$. In particular, we construct a field over which every polynomial in $F[t]$ splits (and is the smallest possible). This field turns out to be an algebraically closed field, i.e., one in which every non-constant polynomial over it splits. We also use the theory developed to answer the classical euclidean construction problems using only straight-edge and compass: the trisection of an angle, doubling of a cube, squaring the circle, and construction of regular $n$-gons, showing how to turn a geometric problem into an algebraic one. The solution of the first two will be done in this chapter as well as the last (assuming a result proven in the next chapter) and the third will be done in §71. Finally, we investigate the difficulty arising from the case of positive characteristic of a field, where irreducible polynomials can have multiple roots over some bigger field.

47. Algebraic Elements

In this section, given a field $F$, we are mostly interested in roots of polynomials with coefficients in $F$. We have seen before that if $f$ is a nonzero polynomial in $F[t]$, then there exists a field $K$ containing $F$
such that \( f \) has a root in \( K \). In fact, we explicitly constructed such a field. We shall investigate this situation in depth.

Let \( K \) be a field containing \( F \) as a subfield, so \( F \subset K \). We shall write this as \( K/F \). We call \( F \) the base field of \( K/F \) and \( K \) the extension field of \( K/F \). If \( E \) is a field with \( F \subset E \subset K \) (as subfields), written \( K/E/F \), we call \( E \) an intermediate field of \( K/F \). We allow \( E \) to be \( K \) or \( F \). A sequence of fields \( F = F_1 \subset \cdots \subset F_n \subset \cdots \) is called a tower of fields. It may be infinite.

**Definition 47.1.** Let \( K/F \) be an extension of fields. Then \( K \) is a vector space over \( F \) by restriction of scalars. We denote the dimension, \( \dim_F K \), of \( K \) as an \( F \)-vector space by \( [K:F] \) and call it the degree of \( K/F \). We call \( K/F \) a finite extension if \( [K:F] \) is finite, and an infinite extension otherwise.

**Examples 47.2.**
1. \( F/F \) is a finite extension of degree one.
2. \( \mathbb{C}/\mathbb{R} \) is a finite extension of degree two.
3. \( \mathbb{R}/\mathbb{Q} \) is an infinite extension.

**Proposition 47.3.** Let \( L/K/F \) be an extension of fields. Then \( L/F \) is a finite extension if and only if both \( L/K \) and \( K/F \) are finite field extensions. Moreover, if \( L/F \) is a finite extension, then
\[
[L:F] = [L:K][K:F].
\]

**Proof.** We prove a stronger statement. Let
\[
\mathcal{B} = \{x_i\}_I \quad \text{be an } F\text{-basis for } K,
\]
\[
\mathcal{C} = \{y_j\}_J \quad \text{be an } K\text{-basis for } L, \text{ and}
\]
\[
\mathcal{D} = \{x_iy_j\}_{I \times J}.
\]

Note \( x_iy_j = x_{i'}y_{j'} \) if and only if \( x_i = x_{i'} \) and \( y_j = y_{j'} \) as \( \mathcal{C} \) is a \( K \)-basis of \( L \) and \( x_i, x_{i'} \) lie in \( K \), hence \( |\mathcal{D}| = |\mathcal{B}||\mathcal{C}| = |I \times J| \).

**Claim.** \( \mathcal{D} \) is an \( F \)-basis for \( L \) and \( |\mathcal{D}| = |I \times J| \).

If we show this, then we are done.

\( \mathcal{D} \) \( F \)-spans \( L \): Let \( z \in L \). Then \( z = \sum_{I,J} \alpha_{ij}y_j \) for some \( \alpha_{ij} \in K \) almost all zero, as \( \mathcal{C} \) spans \( L \) as a vector space over \( K \). As \( \mathcal{B} \) spans \( K \) as a vector space over \( F \), each \( \alpha_{ij} = \sum_I c_{ij}x_i \) for some \( c_{ij} \in F \) almost all zero. Hence we have \( z = \sum_I \sum_J c_{ij}x_iy_j \) lies in the \( F \)-span of \( \mathcal{D} \).

\( \mathcal{D} \) is \( F \)-linearly independent: Suppose that \( \sum_{I \times J} c_{ij}x_iy_j = 0 \) for some \( c_{ij} \) in \( F \), almost all zero. Then \( \sum_J (\sum_I c_{ij}x_i)y_j = 0 \). As \( \mathcal{C} \) is \( K \)-linearly independent, \( \sum_I c_{ij}x_i = 0 \) for all \( j \). As \( \mathcal{B} \) is \( F \)-linearly independent \( c_{ij} = 0 \) for all \( i \) and \( j \). \( \square \)
This field theoretic analogue of Lagrange’s Theorem has two analogous consequences.

**Corollary 47.4.** Let \( L/K/F \) be an extension of fields with \( L/F \) finite. Then \( [K : F] \mid [L : F] \) and \( [L : K] \mid [L : F] \). In particular, if \( [L : F] \) is prime, then \( K = L \) or \( K = F \).

If in the above \( [L : F] \) is prime, we say that there exist no nontrivial intermediate fields in \( L/F \).

**Corollary 47.5.** If \( F_1 \subset \cdots \subset F_n \) is a tower of fields with \( F_n/F_1 \) finite, then
\[
[F_n : F_1] = [F_n : F_{n-1}] \cdots [F_2 : F_1].
\]

Given a field extension, we want to look at intermediate fields.

**Lemma 47.6.** Let \( K \) be a field and \( S \) a non-empty subset of \( K \). Then there exists a unique minimal subfield \( F_0 \) of \( K \) containing \( S \), i.e., if \( S \subset F \subset K \) with \( F \) a field, then \( F_0 \subset F \).

**Proof.** Let
\[
\mathcal{F} = \{ F \mid F \text{ a field satisfying } S \subset F \subset K \}.
\]
As \( K \in \mathcal{F} \), the set \( \mathcal{F} \) is not empty. Set \( F_0 = \bigcap \mathcal{F} \). Certainly, \( S \subset F_0 \). We show that \( F_0 \) works. We first prove that \( F_0 \) is a field. Indeed if \( x, y \in F \) with \( x \neq 0 \), then \( x \) and \( y \) lie in any \( F \in \mathcal{F} \), hence \( x \pm y, xy, x^{-1} \) also lie in any \( F \) in \( \mathcal{F} \), establishing \( F_0 \) is a field. To finish, we show that \( F_0 \) is the unique minimal field containing \( F_0 \). But this is clear, for if \( E \in \mathcal{F} \), then \( F_0 \subset \bigcap \mathcal{F} \subset E \).

**Notation 47.7.** Let \( K/F \) be a field extension and \( X \) a non-empty subset of \( K \). Let \( S = F \cup X \). The field \( F_0 \) in Lemma 47.6 will be denoted by \( F(X) \). If \( X = \{ \alpha_1, \ldots , \alpha_n \} \), we shall write \( F(\alpha_1, \ldots , \alpha_n) \) for \( F(X) \). In particular, if \( \alpha \in K \), then
\[
F(\alpha) = \text{the unique minimal intermediate field of } K/F \text{ containing } \alpha
\]
and will be called a simple (or primitive) extension of \( F \).

**Remark 47.8.** Let \( K/F \) be field extension and \( X \) a non-empty subset of \( K \). As usual \( F[X] \) will denote the ‘polynomials’ in the \( x \in X \) with coefficients in \( F \), i.e., if we let \( t_x \), with \( x \in X \), be independent variables and \( T = \{ t_x \mid x \in X \} \), then \( F[X] \) is the homomorphic image of the evaluation map \( e : F[T] \rightarrow K \) given by \( t_x \mapsto x \) for all \( x \) in \( X \) (where \( F[T] = F[t_x]_{x \in X} \)). As \( F[X] \subset K \), the ring \( F[X] \) is a domain. (So \( \ker e \) is a prime ideal in \( F[T] \).) By the universal property of quotient fields (Theorem 26.14), we must have \( F \subset F[X] \subset qf(F[X]) \subset K \) and
Examples 47.9. 1. \( \mathbb{R}[\sqrt{-1}] = \mathbb{R}(-1) = \mathbb{C} \), as all are 2-dimensional real vector spaces.

2. Let \( d \in \mathbb{Z} \). Then

\[
\mathbb{Q}[\sqrt{d}] = \mathbb{Q}(\sqrt{d}) = qf(\mathbb{Z}[\sqrt{d}])
\]

a vector space over \( \mathbb{Q} \) of dimension at most two and exactly two if \( d \notin \{x^2 \mid x \in \mathbb{Z}\} \) in which case \( \{1, \sqrt{d}\} \) is a \( \mathbb{Q} \)-basis.

3. If \( F \) is a field, then we have

\[
F < F[t] < F(t) := qf(F[t]), \quad F[t]^\times = F^\times, \quad \text{and} \quad F(t)^\times = F(t) \setminus \{0\}.
\]

The last follows as \( F(t) = \{f/g \mid f, g \in F[t] \text{ with } g \neq 0\} \) is a field.

Definition 47.10. Let \( K/F \) be a field extension and \( \alpha \) an element of \( K \). We say

1. \( \alpha \) is algebraic over \( F \) if there exists a nonzero polynomial \( f \) in \( F[t] \) satisfying \( f(\alpha) = 0 \), i.e., \( \alpha \) is a root of a nonzero polynomial with coefficients in \( F \) (hence of a monic polynomial in \( F[t] \)).

2. \( \alpha \) is called a transcendental element over \( F \) if it is not algebraic over \( F \).

Examples 47.11. Let \( F \) be a field.

1. Every element \( \alpha \) in \( F \) is algebraic over \( F \) as \( \alpha \) is a root of the polynomial \( t - \alpha \) in \( F[t] \).

2. Let \( a, b \) be elements in \( \mathbb{Q} \) and \( \alpha = a + b\sqrt{2} \in \mathbb{C} \). Then \( \alpha \) is a root of \( t^2 - 2at + (a^2 + b^2) \) in \( \mathbb{Q}[t] \), so algebraic over \( \mathbb{Q} \).

3. Suppose \( \mathbb{Q} \subset F \subset \mathbb{C} \). [It suffices to say \( F \subset \mathbb{C} \) – why?] and \( d \in \mathbb{Q} \). Then \( F[\sqrt{d}] = F(\sqrt{d}) \), e.g., \( F = \mathbb{Q} \) or \( \mathbb{R} \) Let \( \alpha \) be an element in \( F[\sqrt{d}] \). We can write \( \alpha = a + b\sqrt{d} \) for some \( a, b \in F \) (why?). Then \( \alpha \) is a root of the polynomial \( t^2 - 2at + (a^2 - b^2d) \) in \( F[t] \), hence is algebraic over \( F \):

If \( \sqrt{d} \) lies in \( F \), then \( F = F[\sqrt{d}] = F(\sqrt{d}) \), so suppose not, i.e., suppose that \( d \) is not a square in \( F \). If \( \alpha = a + b\sqrt{d} \) and \( 0 \neq \beta = c + e\sqrt{d} \) with \( a, b, c, e \in F \), then

\[
\frac{\alpha}{\beta} = \frac{a + b\sqrt{d}}{c + e\sqrt{d}} = \frac{(a + b\sqrt{d})(c + e\sqrt{d})}{c^2 - e^2d}.
\]
4. Lindemmann’s Theorem says that the element \( \pi \) is transcendental over \( \mathbb{Q} \). This is not easy \([\text{cf. Section 71}]\). It is easier to show that the real number \( e \) is transcendental over \( \mathbb{Q} \) \([\text{cf. Section 69}]\), as is the Liouville number \( \sum_{i=0}^{10^{10}} 1/(i!) \) \([\text{cf. Section 68}]\). As \( \mathbb{R} \) is uncountable and the set of roots of all nonzero polynomials over \( \mathbb{Q} \) is countable, since there are countably many polynomials over \( \mathbb{Q} \) each having finitely many roots, “most” elements in \( \mathbb{R} \) are transcendental. However, it is usually very difficult to show that a specific real number is transcendental over \( \mathbb{Q} \). It is still unknown if \( \pi^\pi \) is transcendental over \( \mathbb{Q} \).

5. The real numbers \( \pi \) and \( e \) are algebraic over the field \( \mathbb{R} \).

6. The complex number \( e^{\sqrt{-1}\pi} \) is algebraic over \( \mathbb{Q} \) as it is equal to \(-1\), but it can be shown that the number \( e^\pi \) is transcendental over \( \mathbb{Q} \).

7. Let \( L/K/F \) be field extensions. If \( \alpha \) is an element in \( L \) algebraic over \( F \) then it is algebraic over \( K \), since \( F[t] \subset K[t] \). However, in general, if \( \alpha \) is algebraic over \( K \) it need not be algebraic over \( F \).

8. Let \( K/F \) be a field extension and \( u \) an element in \( K \) that is transcendental over \( F \). We know that the evaluation map \( e_u : F[t] \to F[u] \) given by \( f \mapsto f(u) \) is a ring epimorphism and an \( F \)-linear transformation. Suppose that \( f \) lies in \( \ker e_u \). As \( u \) is transcendental over \( F \), we must have \( f = 0 \), hence \( e_u \) is a ring monomorphism. As it is also surjective, \( e_u : F[t] \to F[u] \) is an isomorphism. The variable \( t \) is not a unit in \( F[t] \), since \( F[t]^\times = F^\times \). Therefore, we also have \( F[t] \not< F(t) \), so \( F[u] \not< F(F[u]) = F(u) \). In addition, as \( \{t^i \mid i \geq 0\} \) is an \( F \)-basis for \( F[t] \), we have \( \{u^i \mid i \geq 0\} \) is an \( F \)-basis for \( F[u] \). In particular, \( F[u] \) and \( F(u) \) are both infinite dimensional vector spaces over \( F \).
Note: By (4), this means that $\mathbb{Q}[\pi] \cong \mathbb{Q}[t]$. But $\mathbb{R}[\pi] = \mathbb{R} < \mathbb{R}[t]$, so being transcendental is a relative notion depending on the base field $F$.

The main theorem about an algebraic element over a field is given by the following foundational result, much of which we have implicitly seen before.

**Theorem 47.12.** Let $K/F$ be an extension of fields and $\alpha$ an element in $K$ that is algebraic over $F$. Then there exists a unique monic irreducible polynomial $m_F(\alpha)$ in $F[t]$ satisfying all of the following:

1. $\alpha$ is a root of $m_F(\alpha)$.
2. If $\alpha$ is a root of a nonzero polynomial $g$ in $F[t]$, then $m_F(\alpha) | g$ in $F[t]$ and $\deg g \geq \deg m_F(\alpha)$ with equality if and only if $g = am_F(\alpha)$ for some nonzero element $a \in F$.
3. Let $n = \deg m_F(\alpha)$. Then $F(\alpha) = F[\alpha]$ and $\{1, \alpha, \ldots, \alpha^{n-1}\}$ is an $F$-basis for $F[\alpha] = F(\alpha)$. In particular,

$$[F(\alpha) : F] = \dim_F F[\alpha] = \deg m_F(\alpha).$$

The polynomial $m_F(\alpha)$ in the theorem is called the **minimal or irreducible polynomial** of $\alpha$ over $F$. The integer $\deg m_F(\alpha) = [F(\alpha) : F]$ is called the **degree of $\alpha$ over $F$**.

**Proof.** We begin the proof by modifying the argument in Construction 44.1. Let $e_\alpha : F[t] \to F[\alpha]$ be given by $f \mapsto f(\alpha)$, evaluation at $\alpha$, a ring epimorphism and $F$-linear transformation of vector spaces over $F$. As $\alpha$ is algebraic over $F$, $\alpha$ is the root of a nonzero polynomial in $F[t]$ which may be assumed to be monic. Therefore, there exists a positive integer $N$ such that $\alpha^N$ lies in the $F$-vector space $\sum_{i=0}^{N-1} F\alpha^i$. It follows by induction that $\alpha^{N+j}$ lies in $\sum_{i=0}^{N-1} F\alpha^i$ for all $j \geq 0$, so $F[\alpha] = \sum_{i=0}^{N-1} F\alpha^i$ is a finite dimensional vector space over $F$. Since $F[t]$ is an infinite dimensional vector space over $F$, it follows that $e_\alpha$ cannot be a monomorphism, so $0 < \ker e_\alpha < F[t]$ ($e_\alpha(1) = 1$, so $e_\alpha$ is not the trivial map).

Since $F[t]$ is euclidean, it is a PID, hence there exists a unique non-constant monic polynomial $f$ in $F[t]$ such that $\ker e_\alpha = (f)$. [It is unique, as generators of $\ker e_\alpha$ only differ by a nonzero element in $F$.] By the First Isomorphism Theorem, $F[\alpha] = \text{im } e_\alpha \cong F[t]/(f)$. Let $n = \deg f$.

**Claim.** $f = m_F(\alpha)$ works.

As $F[t]/(f) \cong F[\alpha] \subset K$, it is a domain, so $0 < (f) < F[t]$ is a prime ideal. In particular, $f$ must be an irreducible polynomial. Since $F[t]$ is
a PID, it must be a maximal ideal, hence \( F[t]/(f) \cong F[\alpha] \) is a field, so \( F[\alpha] = F(\alpha) \).

We next show that \( \mathcal{B} = \{1, \alpha, \ldots, \alpha^{n-1}\} \) is an \( F \)-basis for \( F[\alpha] \). Let \( g \in F[t] \). As \( F[t] \) is euclidean (under the degree map), \( g = fq + r \) for some \( q, r \in F[t] \) with \( r = 0 \) or \( \deg r < \deg f \). Since \( e_\alpha \) is a ring homomorphism, \( g(\alpha) = f(\alpha)q(\alpha) + r(\alpha) = r(\alpha) \), which lies in \( \sum_{i=0}^{n-1} F\alpha^i \), the span of \( \mathcal{B} \). Moreover, \( g(\alpha) = 0 \) if and only if \( r(\alpha) = 0 \). As \( r = 0 \) or \( \deg r < \deg f \), we have \( r(\alpha) = 0 \) if and only if \( r \in \ker e_\alpha \) if and only if \( f \mid r \) in \( F[t] \) if and only if \( r = 0 \). We conclude that if \( g(\alpha) = \sum_{i=0}^{n-1} a_i\alpha^i = 0 \) with \( a_i \in F \), then each \( a_i = 0 \). Consequently, \( \mathcal{B} \) is also \( F \)-linearly independent, hence an \( F \)-basis for \( F[\alpha] \). The other needed properties of \( f \) now follow easily.

**Remarks 47.13.** Let \( K/E/F \) be field extensions and \( \alpha \in K \) an algebraic element over \( F \).

1. \( m_E(\alpha) \mid m_F(\alpha) \) in \( E[t] \) and \( \deg m_E(\alpha) \leq \deg m_F(\alpha) \) [Of course, they can be different.]

2. If \( f \in F[t] \) is monic and irreducible with \( f(\alpha) = 0 \), then \( f = m_F(\alpha) \).

3. If \( g \in F[t] \) satisfies \( \deg g = [F(\alpha) : F] \) and \( g(\alpha) = 0 \), then \( g \) is irreducible.

**Proposition 47.14.** (Characterization of algebraic elements) Let \( K/F \) be a field extension and \( \alpha \in K \). Then \( \alpha \) is algebraic over \( F \) if and only if \( [F(\alpha) : F] \) is finite.

**Proof.** (\( \Rightarrow \)): This follows from the Theorem.

(\( \Leftarrow \)): Suppose that \( n = [F(\alpha) : F] \). Then \( \{1, \alpha, \ldots, \alpha^n\} \) is \( F \)-linearly dependent. Therefore, there exist \( a_i \in F, 1 \leq i \leq n \), not all zero such that \( \sum_{i=0}^{n} a_i\alpha^i = 0 \), so \( \alpha \) is a root of the nonzero polynomial \( \sum_{i=0}^{n} a_i t^i \) in \( F[t] \).

**Corollary 47.15.** Let \( K/F \) be a finite field extension. Then every element in \( K \) is algebraic over \( F \).

**Proposition 47.16.** (Characterization of transcendental elements) Let \( K/F \) be a field extension and \( u \in K \). Then the following are equivalent:

1. \( u \) is transcendental over \( F \).
2. \( F[u] \cong F[t] \).
3. \( F[u] < F(u) \).
4. \( [F(u) : F] \) is infinite.
5. \( \dim_F F(u) \) is infinite.
Corollary 47.17. Let $K/F$ be a field extension. If $a_1,\ldots,a_r$ are elements of $K$ algebraic over $F$, then $F[a_1,\ldots,a_r] = F(a_1,\ldots,a_r)$ and $[F(a_1,\ldots,a_r) : F] \leq \prod_{i=1}^{r}[F(a_i) : F]$ (with strict inequality possible).

Proof. As $a_r$ is algebraic over $F$, it is algebraic over $F[a_1,\ldots,a_{r-1}]$. By induction and the $r=1$ case (that we have done), $F[a_1,\ldots,a_r] = F[a_1,\ldots,a_{r-1}][a_r] = F(a_1,\ldots,a_{r-1})[a_r] = F(a_1,\ldots,a_r)$ and

$$[F(a_1,\ldots,a_r) : F] = \deg m_{F(a_1,\ldots,a_{r-1})}(a_r) \prod_{i=1}^{r-1}[F(a_i) : F] \leq \prod_{i=1}^{r}[F(a_i) : F].$$

Corollary 47.18. Let $A$ be a domain containing a field $F$. If $A$ is a finite dimensional vector space over $F$, then $A$ is a field.

Proof. If $A$ has an $F$-basis $\mathcal{B}$, then $A = F[\mathcal{B}]$ and its quotient field $K = F(\mathcal{B})$ is a finite dimensional vector space over $F$ on basis $\mathcal{B}$, hence each element of $\mathcal{B}$ is algebraic over $F$. It follows that $A = K$ by the corollary.

If $K/F$ is a field extension, we arrive at the important conclusion that the set of elements in $K$ algebraic over $F$ form an intermediate field.

Theorem 47.19. Let $K/F$ be a field extension with $a, b$ elements in $K$ algebraic over $F$. Then $a \pm b, ab$ and $b^{-1}$, if $b \neq 0$, are algebraic over $F$. In particular, the set of elements in $K$ algebraic over $F$ forms an intermediate field of $K/F$.

Proof. Let $\gamma = a \pm b, ab$ or $b^{-1}$, if $b \neq 0$. Then $F \subset F(\gamma) \subset F(a, b)$, so

$$[F(\gamma) : F] \leq [F(a, b) : F] \leq [F(a) : F][F(b) : F] < \infty.$$
**Definition 47.20.** Let $K/F$ be a field extension. We say that it is an *algebraic extension* (or simply *algebraic*) if every element of $K$ is algebraic over $F$. If $K/F$ is not algebraic, we say that it is a *transcendental extension*.

We have shown that every finite field extension is algebraic. The converse is false, as we shall see below. First, however, we show that being an algebraic extension is transitive. In fact, we have

**Theorem 47.21.** Let $L/K/F$ be field extensions. Then $L/F$ is algebraic if and only if both $L/K$ and $K/F$ are algebraic.

**Proof.** ($\Rightarrow$): Any element in $L$ algebraic over $F$ is clearly algebraic over $K$ and any element in $K$ lies in $L$ so is algebraic over $F$.

($\Leftarrow$): Let $\alpha$ be an element of $L$. We must show that $\alpha$ is algebraic over $F$. As it is algebraic over $K$, the minimal polynomial $m_K(\alpha) = \sum_{i=0}^{n} c_i t^i$ in $K[t]$ exists. Let $E = F(c_0, \ldots, c_{n-1}) \subset K$. Since the elements $c_0, \ldots, c_{n-1}$ are algebraic over $F$, we have $[E : F] < \infty$. As $m_K(\alpha) \in E[t]$, we also have $[E(\alpha) : E] < \infty$, so

$$[F(\alpha) : F] \leq [E(\alpha) : F] = [E(\alpha) : E][E : F] < \infty$$

and $\alpha$ is algebraic over $F$. \hfill $\Box$

The above trick is very useful. If you want to prove a result and can reduce to a finite amount of data, you can often reduce to checking in a nicer set, e.g., in the above we reduced from an algebraic extension to a finite extension, because we needed to attach only finitely many elements algebraic over the base field. In commutative ring theory, a finite amount of data over a commutative ring, often reduces a proof to a Noetherian ring containing it by the Hilbert Basis Theorem.

**Examples 47.22.** 1. $\mathbb{C}/\mathbb{R}$ and $\mathbb{Q}(\sqrt{d})/\mathbb{Q}$, with $d$ an integer, are algebraic extensions.

2. $\mathbb{R}/\mathbb{Q}$ is transcendental as $\mathbb{R}$ is uncountable (by Cantor’s Theorem) and $\mathbb{Q}$ is countable.

3. Let 

$$\Omega := \{\alpha \in \mathbb{C} \mid \alpha \text{ algebraic over } \mathbb{Q}\}.$$ 

$\Omega$ is a countable field as it is the set of roots of countably many nonzero polynomials, each having finitely many roots.

**Claim.** $\Omega/\mathbb{Q}$ is not finite:

Let $p$ be a (positive) prime in $\mathbb{Z}$ and $n$ a positive integer. By Eisenstein’s Criterion (34.11(1)), $t^n - p$ in $\mathbb{Q}[t]$ is irreducible and has a
Let \( \alpha \) be a root in \( \Omega \). In particular, \(|\Omega : \mathbb{Q}| \geq |\mathbb{Q}(\alpha) : \mathbb{Q}| = \deg m_\mathbb{Q}(\alpha) = \deg(t^n-p) = n \), establishing the claim.

4. \( \Omega \) in (3) is algebraically closed, i.e., every non-constant polynomial in \( \Omega[t] \) has a root in \( \Omega \). To see this, we use the Fundamental Theorem of Algebra (to be proved in 56.12 below) that \( \mathbb{C} \) is algebraically closed. If \( f \) is a non-constant polynomial in \( \Omega[t] \), it has a root \( \alpha \) in \( \mathbb{C} \), so \( \Omega(\alpha)/\Omega \) is algebraic. As \( \Omega/\mathbb{Q} \) is algebraic, it follows that so is \( \Omega(\alpha)/\mathbb{Q} \), hence \( \alpha \) is algebraic over \( \mathbb{Q} \) hence is in \( \Omega \). As \( \mathbb{C} \) is uncountable, we must have \( \Omega < \mathbb{C} \).

5. The same argument as in (4) shows if \( K/F \) with \( K \) algebraically closed and \( \Omega \) the set of elements of \( K \) that are algebraic over \( F \), then \( \Omega \) is an algebraically closed field. Moreover, by definition, if \( F < E < \Omega \) is a field, then \( E \) is not algebraically closed. We call \( \Omega \) an algebraic closure of \( F \). We shall see in \( \S 50 \), that all such algebraic closures of a field \( F \) are unique, up to a field isomorphism fixing \( F \).

**Calculation 47.23.** 1. Let \( p \in \mathbb{Z}^+ \) be a prime. Then the polynomial \( f = t^{p-1} + t^{p-2} + \ldots + t + 1 \) is irreducible in \( \mathbb{Q}[t] \) by the application of Eisenstein’s Criterion in 34.11(3). Let \( \zeta_p \) be a root of \( f \) in \( \mathbb{C} \). Then \( f = m_\mathbb{Q}(\zeta_p) \), so \([\mathbb{Q}(\zeta_p) : \mathbb{Q}] = \deg f = p - 1 = \varphi(p) \) (where \( \varphi \) is the Euler phi-function), and \( \{1, \zeta_p, \ldots, \zeta_p^{p-2}\} \) is a \( \mathbb{Q} \)-basis for \( \mathbb{Q}(\zeta_p) \).

[Note that \( t^p - 1 = (t - 1)(t^{p-1} + t^{p-2} + \cdots + t + 1) \) is a factorization into irreducibles in \( \mathbb{Q}[t] \).

2. Let \( \alpha \) in \( \mathbb{C} \) be a root of the irreducible polynomial \( f = t^3 - 2t + 2 \) in \( \mathbb{Q}[t] \) (irreducible by Eisenstein’s Criterion). So \([\mathbb{Q}(\alpha) : \mathbb{Q}] = \deg f = \deg m_\mathbb{Q}(\alpha) = 3 \) and \( B = \{1, \alpha, \alpha^2\} \) is a \( \mathbb{Q} \)-basis for \( \mathbb{Q}(\alpha) \). As \( \alpha^3 = 2\alpha - 2 \) in \( \mathbb{Q}(\alpha) \), we have \(-\alpha^{-1} = (\alpha^2/2) - 1 \) in \( \mathbb{Q}(\alpha) \). Let \( \beta \in \mathbb{Q}(\alpha) \), so \( \mathbb{Q}(\beta) \subset \mathbb{Q}(\alpha) \), hence \([\mathbb{Q}(\beta) : \mathbb{Q}] \mid [\mathbb{Q}(\alpha) : \mathbb{Q}] = 3 \). Therefore, either \( \beta \in \mathbb{Q} \) or \( \mathbb{Q}(\beta) = \mathbb{Q}(\alpha) \). For example, suppose that \( \beta = \alpha^2 - \alpha \). As \( B \) is \( \mathbb{Q} \)-linearly independent, \( \beta \notin \mathbb{Q} \), so \( \mathbb{Q}(\beta) = \mathbb{Q}(\alpha) \) and \( C = \{1, \beta, \beta^2\} \) is a \( \mathbb{Q} \)-basis for \( \mathbb{Q}(\beta) \) and \( \{1, \beta, \beta^2, \beta^3\} \) is \( \mathbb{Q} \)-linearly dependent. Computation shows that

\[
1 = 1, \quad \beta = \alpha^2 - \alpha, \quad \beta^2 = 3\alpha^2 - 6\alpha + 4
\]

and

\[
(*) \quad \beta^3 = 16\alpha^2 - 28\alpha + 18.
\]

It follows that the change of basis matrix for the basis \( C \) to the basis \( B \) is

\[
[Id]_{C,B} = \begin{pmatrix} 1 & 0 & 4 \\ 0 & -1 & -6 \\ 0 & 1 & 3 \end{pmatrix}.
\]
Inverting this matrix to \([Id]_{B,C}\) then shows that
\[
1 = 1, \quad \alpha = -\frac{\beta^2}{3} + \beta + \frac{4}{3}, \quad \alpha^2 = -\frac{\beta^2}{3} + 2\beta + 4 + \frac{4}{3},
\]
and substituting this into \((*)\) yields \(\beta^3 - 4\beta^2 - 4\beta - 2 = 0\), i.e., \(\beta\) is a root of monic \(g = t^3 - 4t^2 - 4t - 2\) in \(\mathbb{Q}[t]\). This must be the irreducible polynomial \(m_{\mathbb{Q}(\beta)}\) as \(g\) has the same degree and satisfies \(g(\beta) = 0\). [Of course, we also know that it is irreducible by Eisenstein’s Criterion.]

3. It is easy to check that \([\mathbb{Q}(\sqrt{2}, \sqrt{3}) : \mathbb{Q}] = 4\) with \(\{1, \sqrt{2}, \sqrt{3}, \sqrt{6}\}\) a \(\mathbb{Q}\)-basis for \(\mathbb{Q}(\sqrt{2}, \sqrt{3})\) (by showing that \(\sqrt{3}\) cannot lie in \(\mathbb{Q}(\sqrt{2})\)). [Cf. Exercise 49.18(12) for the generalization.] Using this we determine \([\mathbb{Q}(\sqrt{2}, \sqrt{3}) : \mathbb{Q}] = 1, 2\) or 4. Let \(\alpha = \sqrt{2} + \sqrt{3}\). Then
\[
1 = 1, \quad \alpha = \sqrt{2} + \sqrt{3}, \quad \alpha^2 = 5 + 2\sqrt{6}.
\]
These must be \(\mathbb{Q}\)-linearly independent (why?). It follows that the set \(\{1, \alpha, \alpha^2, \alpha^3\}\) must be a \(\mathbb{Q}\)-basis for \(K\). Now \(\alpha^2 = 5 + 2\sqrt{6}\), so the polynomial \(g = t^2 - (5 + 2\sqrt{6})\) lying in \(\mathbb{Q}(\sqrt{6})[t]\) has a root \(\alpha\) in \(K\) with \(\alpha \not\in \mathbb{Q}(\sqrt{6})\). Consequently, \([K : \mathbb{Q}(\sqrt{6})] = 2\) and \(g = m_{\mathbb{Q}(\sqrt{6})}(\alpha)\). Computation shows that
\[
h = (t^2 - (5 + 2\sqrt{6}))(t^2 - (5 - 2\sqrt{6})) = t^4 - 10t^2 + 1,
\]
so it lies in \(\mathbb{Q}[t]\) and has \(\alpha\) as a root in \(K\) with \(\deg h = \deg m_{\mathbb{Q}(\sqrt{6})}(\alpha) = 4\). Thus \(h = m_{\mathbb{Q}}(\alpha)\).

We shall show later that if \(F\) is a field of characteristic zero and \(K/F\) a finite extension, that there always exists a \(\theta\) in \(K\) such that \(K = F(\theta)\). [Cf. Theorem 56.9.] In general, such a \(\theta\) is hard to find, although many different ones work.

**Exercises 47.24.**

1. Let \(R\) be a domain containing a field \(F\) and over which it is a finite dimensional vector space. Show that \(R\) is a field. In particular, show if \(D\) is a division ring containing \(F\) and \(\alpha \in D\) is a root of a nonzero polynomial in \(F[t]\), then \(F[\alpha]\) is a field. [Cf. A finite division ring is a field.]

2. (a) Find \(u \in \mathbb{R}\) such that \(\mathbb{Q}(u) = \mathbb{Q}(\sqrt{2}, \sqrt{5})\).
(b) Describe how you how you would find all \(w \in \mathbb{Q}(\sqrt{2}, \sqrt{5})\) such that \(\mathbb{Q}(w) = \mathbb{Q}(\sqrt{2}, \sqrt{5})\).
3. If $a, b \in K$ are algebraic over $F$ are of degree $m, n$ respectively with $m, n$ relatively prime, show that $[F(a, b) : F] = mn$.

4. Show that $[\mathbb{Q}(\sqrt{2}, \sqrt{3}) : \mathbb{Q}] = 4$ with $\{1, \sqrt{2}, \sqrt{3}, \sqrt{6}\}$ a $\mathbb{Q}$-basis for $\mathbb{Q}(\sqrt{2}, \sqrt{3})$.

5. If $|F| = q < \infty$ show:
   (a) There exists a prime $p$ such that $\text{char } F = p$.
   (b) $q = p^n$ some $n$.
   (c) $a^q = a$ for all $a \in F$.
   (d) If $b \in K$ is algebraic over $F$ then $b^{q^m} = b$ for some $m > 0$.

6. Let $u$ be a root of $f = t^3 - t^2 + t + 2$ in $\mathbb{Q}[t]$ and $K = \mathbb{Q}(u)$.
   (a) Show that $f = m\mathbb{Q}(u)$.
   (b) Express $(u^2 + u + 1)(u^2 - u)$, and $(u-1)^{-1}$ in the form $au^2 + bu + c$, for some $a, b, c \in \mathbb{Q}$.

7. Let $\zeta = \cos \frac{\pi}{6} + \sqrt{-1} \sin \frac{\pi}{6}$ in $\mathbb{C}$. Show that $\zeta^{12} = 1$ but $\zeta^r \neq 1$ for $1 \leq r < 12$. Show also that $[\mathbb{Q}(\zeta) : \mathbb{Q}] = 4$ and find $m_{\mathbb{Q}}(\zeta)$.

8. Let $K = F(u)$, $u$ algebraic over $F$ and degree of $u$ odd. Show that $K = F(u^2)$.

9. Let $u$ be transcendental over $F$ and $F < E \subset F(u)$. Show that $u$ is algebraic over $E$.

10. Let $f, g \in F[t]$ be relatively prime and suppose that $u = f/g$ lies in $F(t) \setminus F$. Show that $F(t)/F(u)$ is finite of degree $d = \max\{\deg(f), \deg(g)\}$.

11. If $f = t^n - a \in F[t]$ is irreducible and $u \in K$ is a root of $f$ and $n/m \in \mathbb{Z}$, show that $[F(u^m) : F] = \frac{n}{m}$. What is $m_{F}(u^m)$?

12. If $a^n$ is algebraic over a field $F$ for some $n > 0$, show that $a$ is algebraic over $F$.

48. Addendum: Transcendental Extensions

What can we say about a field extension of a field that is not algebraic, i.e., is transcendental? In this short section, we answer this question, leaving most of the details to the reader.

**Definition 48.1.** Let $K/F$ be an extension of fields and suppose that $\mathcal{B} = \{x_i \mid i \in I\}$ a subset of $K$. We say that the set $\mathcal{B}$ is algebraically independent over $F$ if $x_i \neq x_j$ for $i \neq j$ and if $f \in F[\mathcal{B}]$ satisfies $f(x_i) = 0$, then $f = 0$. This is equivalent to

$$\{x_1^{n_{11}} \cdots x_r^{n_{1r}} \mid n_{ij} \geq 0 \text{ for all } i_j \text{ in } I, \text{ some } r \geq 0\}$$

is an $F$-linearly independent set for the $F$-vector space $F[\mathcal{B}]$. 
For example, a set of indeterminants $T = \{t_i \mid i \in I\}$ is an algebraically independent set in the quotient field $F(T)$ of $F[T]$. A maximal algebraically independent subset $\mathcal{B}$ in an extension field $K$ of $F$ is called a transcendence basis for $K$ over $F$.

Remarks 48.2. Let $K/F$ be an extension of fields.
1. There exists a transcendence basis for $K$ over $F$. If $K/F$ is algebraic, it is the empty set. A proof of this is analogous to the proof of the existence of bases for vector spaces.
2. Is $\mathcal{B}$ is a transcendence basis for $K$ over $F$ then $K/F(\mathcal{B})$ is algebraic.
3. The set $\{t_1, \ldots, t_n\}$ is a transcendence basis for $F(t_1, \ldots, t_n)$ over $F$.
4. The set $\mathcal{B} = \{x_i\}_{I} \subset K$ is a transcendence basis for $K$ over $F$ if and only if $\mathcal{B}$ is a maximal subset of $K$ satisfying $F(\mathcal{B}) \cong F(\{t_i\}_I)$.
5. The extension $K/F$ is called finitely generated if there exists a finite set $S \subset K$ such that $K = F(S)$. If $K/F$ is finitely generated, then $K$ has a finite transcendence basis over $F$.

The analogue of dimension of vector spaces holds for the cardinality of transcendence bases. The proof for the case that $K/F$ finitely generated is analogous to the so called Replacement Theorem used to prove that dimension is well-defined for finite dimensional vector spaces. We shall only prove this case in some detail.

Proposition 48.3. Let $K/F$ be an extension of fields. Then any two transcendence bases of $K$ over $F$ have the same cardinality.

Proof. Although this follows from the fact that the cardinality of any basis for a fixed vector space is well-defined, we present a proof in the case that $K/F$ is a field extension and $K$ has a finite transcendence basis $\{x_1, \ldots, x_n\}$ over $F$. Suppose that $\{y_1, \ldots, y_m\}$ is a subset of $K$ algebraically independent over $F$. We must show $m \leq n$. As $\{x_1, \ldots, x_n\}$ is a maximal algebraically independent set over $F$, the set $\{y_1, x_2, \ldots, x_n\}$ must be algebraically dependent, i.e., is not algebraically independent. Therefore, there exists a nonzero polynomial $f \in F[t_0, t_1, \ldots, t_n]$ satisfying $f(y_1, x_1, \ldots, x_n) = 0$. By assumption, $y_1$ occurs non-trivially in $f$ as well as some $x_i$, say $x_1$. Then $x_1$ is algebraic over $F(y_1, x_2, \ldots, x_n)$. Hence both $K/F(x_1, \ldots, x_n)$ and $F(y_1, x_2, \ldots, x_n)$ are algebraic over $F(y_1, x_2, \ldots, x_n)$. It follows that $K/F(y_1, x_2, \ldots, x_n)$ is algebraic. A subset $\mathcal{B}$ of $\{y_1, x_2, \ldots, x_n\}$ must be a transcendence basis of $K$ over $F$ and it must include $y_1$, for if $\mathcal{B} < \{y_1, x_2, \ldots, x_n\}$, reversing the argument, would show that a proper subset of $\{x_1, \ldots, x_n\}$ would be a transcendence basis for $K$ over $F$. It follows that $\{y_1, x_2, \ldots, x_n\}$ is a transcendence basis for $K$ over $F$. 
The argument continues, as in the proof of the Replacement Theorem for vector spaces, by throwing in $y_2$ and seeing that we have to throw away some $x_i$ in $\{x_2, \ldots, x_n\}$, etc. It follows that $m \leq n$. □

Let $K/F$ be a field extension. Then the cardinality of a transcendence basis $\mathcal{B}$ for $K$ over $F$ is called the transcendence degree of $K/F$ and denoted by $\text{tr deg}_F K$. It follows that $K/F$ is finitely generated if and only if $\text{tr deg}_F K$ is finite and $K/F(\mathcal{B})$ is a finite extension.

**Proposition 48.4.** Let $L/K/F$ be field extensions. Then $\text{tr deg}_F L$ is finite if and only if both $\text{tr deg}_F K$ and $\text{tr deg}_K L$ are finite. Moreover, if $\text{tr deg}_F L$ is finite, then $\text{tr deg}_F L = \text{tr deg}_K L + \text{tr deg}_F K$.

**Exercises 48.5.**
1. Prove that if $K/F$ is an extension of fields, then there exists a transcendence basis for $K$ over $F$.
2. If $K/F$ is an extension of fields, prove that $\mathcal{B} = \{x_i\}_I \subset K$ is a transcendence basis for $K/F$ if and only if $\mathcal{B}$ is a maximal subset of $K$ satisfying $F(\mathcal{B}) \cong F(\{t_i\}_I)$ and $F[\mathcal{B}] \cong F[\{t_i\}_I]$.

**49. Splitting Fields**

We have shown in Kronecker’s Theorem 33.12 that an irreducible polynomial over a field $F$ has a root in an extension field. Indeed we construct such as follows:

Let $F$ be a field and $f$ an irreducible polynomial in $F[t]$. As $F[t]$ is a PID, the non-trivial prime ideal $(f)$ in $F[t]$ is maximal, so $K = F[t]/(f)$ is a field. Let $\varphi : F[t] \to K$ be the canonical epimorphism given by $g \mapsto g + (f)$. The composition $F \subset F[t] \to K$ is monic, since $F$ is simple, so we can, and do, view this as an inclusion of $F$ into $K$, i.e., we identify $a$ in $F$ and $\overline{a}$ in $K$, so $K/F$ is a field extension. Under this identification,

$$h := \sum_i a_i t^i \mapsto \overline{h} = \sum i \overline{a_i t^i} = \sum a_i \overline{t^i} = h(\overline{t})$$

is evaluation at $\overline{t}$. In particular, $\overline{0} = \overline{f} = f(\overline{t})$ in $K$, so $f \in F[t] \subset K[t]$ has a root in $K$.

Continuing, if $g$ lies in $F[t]$ then, as $F[t]$ is euclidean, $g =fq + r$ in $F[t]$ for some $q,r \in F[t]$ with $r = 0$ or $\deg r < \deg f$. Therefore, $\overline{g} = \overline{fq} + \overline{r} = \overline{r}$. It follows that if $n = \deg f$, then $\mathcal{B} = \{1, \overline{t}, \ldots, \overline{t}^{n-1}\}$ spans the vector space $K$ over $F$. As $\overline{g} = 0$ in $K$ if and only $f \mid g$ in
49. SPLITTING FIELDS

K[t], we have $g = 0$ or $\deg f \leq \deg g$. It follows that $B$ is an $F$-basis for $K$ (why?), $f = (\text{lead } f)m_F(t)$, and $n = \deg f = \deg m_F(t) = [K : F]$.

So we have proved

**Proposition 49.1.** Let $F$ be a field and $f$ an irreducible polynomial in $F[t]$. Then there exists an extension field $K/F$ of degree $\deg f$ such that $f$ has a root in $K$.

Generalizing (as we did before), we have:

**Theorem 49.2.** (Kronecker’s Theorem) Let $F$ be a field and $f$ a non-constant polynomial in $F[t]$. Then there exists an algebraic extension $K/F$ such that $f$ has a root in $K$ and $[K : F] \leq \deg f$.

**Proof.** Let $f_1 | f$ in the UFD $F[t]$ with $f_1$ irreducible. Then $f$ has a root in the field $K = F[t]/(f_1)$, as $f_1$ does, with $[K : F] = \deg f_1 \leq \deg f$. □

**Definition 49.3.** Let $K/F$ be an extension of fields and $f$ a non-constant polynomial in $F[t]$. We say that $f$ splits over $K$ if $f$ factors into a product of linear polynomials in $K[t]$ and that $K$ is a splitting field of $f$ over $F$ or $K/F$ is a splitting field of $f$ if $K$ is a minimal field extension of $F$ such that $f$ splits over $K$, i.e.,

1. $f$ splits over $K$.
2. Either $K = F$ or whenever $K/E/F$ are field extensions with $E < K$, the polynomial $f$ does not split over $E$.

The existence of splitting fields of a non-constant polynomial now will follow by an easy induction using Kronecker’s Theorem.

**Theorem 49.4.** Let $F$ be a field and $f$ a non-constant polynomial in $F[t]$. Then there exists a splitting field $K$ of $f$ over $F$ satisfying $[K : F] \leq (\deg f)!$.

**Proof.** Let $n = \deg f$. By Kronecker’s Theorem, there exists a field extension $K_1/F$ such that $f = (t - \alpha_1)f_1$ in $K_1[t]$ with $f_1 \in K_1[t]$ and $[K_1 : F] \leq n$. By induction, there exists $K_2$ a splitting field of $f_1$ over $K_1$ satisfying $[K_2 : K_1] \leq (\deg f_1)! = (n - 1)!$. Thus $[K_2 : F] = [K_2 : K_1][K_1 : F] \leq n!$ and $f$ splits over $K_2$. By well-ordering, there exists an intermediate field $K_2/K/F$ with $K$ a splitting field of $f$ over $F$ and $[K : F] \leq n!$. [Is $K = K_2$?] □

**Remarks 49.5.** Let $F$ be a field and $f$ a non-constant polynomial in $F[t]$ of degree $n$, $K$ a splitting field of $f$ over $F$, and $\alpha_1, \ldots, \alpha_n$ the roots of $f$ in $K$. [They may not be distinct.]
1. We have $K = F(\alpha_1, \ldots, \alpha_n)$ and the proof of the theorem shows that $[K : F] \leq (\deg f)! = n! = |S_n|$. One often calls the splitting field $K$ of $f$ over $F$ the root field of $f$ over $F$ for this reason.

2. Let $L/F$ be a field extension with both $K$ and $L$ lying in some larger field. Then $L(\alpha_1, \ldots, \alpha_n)$ is a splitting field of $f$ over $L$. Of course, it is possible that $L = L(\alpha_1, \ldots, \alpha_n)$.

3. If $K/L/F$ are field extensions, then $K$ is a splitting field of $f$ over $L$.

4. If $g \mid f$ in $F[t]$ with $g$ a non-constant polynomial, then $g$ splits over $K$, hence there exists an intermediate field $K/E/F$ such that $E$ is a splitting field of $g$ over $F$.

We turn to uniqueness statements for splitting fields. We shall do this in greater generality, so we need some notation.

**Definition 49.6.** Let $R_1 \subset S_1$ and $R_2 \subset S_2$ be commutative rings and $\varphi : R_1 \to R_2$ and $\tau : S_1 \to S_2$ be ring homomorphisms. We say that $\tau$ lifts or extends $\varphi$ if $\tau|_{R_1} = \varphi$, i.e., we have a commutative diagram:

$$
\begin{array}{ccc}
S_1 & \xrightarrow{\tau} & S_2 \\
inc \downarrow & & \downarrow inc \\
R_1 & \xrightarrow{\varphi} & R_2.
\end{array}
$$

where $inc$ is the inclusion. If $R = R_1 = R_2$, $\varphi = 1_R$, and $\tau$ lifts $\varphi$, we call $\tau$ an $R$-algebra homomorphism. If, in addition, $\tau$ is an isomorphism (respectively, monomorphism, epimorphism), then we call $\tau$ an $R$-algebra isomorphism (respectively, $R$-algebra monomorphism, $R$-algebra epimorphism). (Cf. Remark 25.11.) If $S = S_1 = S_2$ and $\tau$ is an automorphism, we call it an $R$-algebra automorphism. As is standard, if $R = F$ is a field, we call such an $F$-algebra homomorphism (respectively, $F$-algebra monomorphism, $F$-algebra epimorphism, $F$-algebra isomorphism, $F$-algebra automorphism) just an $F$-homomorphism (respectively, $F$-monomorphism, $F$-epimorphism, $F$-isomorphism, $F$-automorphism). Of course, an $F$-algebra homomorphism is also a $F$-vector space linear transformation, but the converse is not true in general. Since these standard notations conflict with our usage of a module homomorphism if $F$ is a field, a $F$-vector space homomorphism will be called an $F$-linear homomorphism or an $F$-linear transformation from now on.

**Examples 49.7.** 1. Let $\varphi : R \to S$ be a ring homomorphism of commutative rings. Define $\tilde{\varphi} : R[t] \to S[t]$ by $\sum a_i t^i \mapsto \sum \varphi(a_i) t^i$, a ring homomorphism extending $\varphi$. Given such a $\varphi$, we shall always
let \( \varphi \) denote this map. If \( \varphi \) is a monomorphism, respectively an epimorphism, isomorphism, then \( \tilde{\varphi} \) is a monomorphism, respectively an epimorphism, isomorphism.

2. Complex conjugation \( - : \mathbb{C} \to \mathbb{C} \) given by \( a + b\sqrt{-1} \mapsto a - b\sqrt{-1} \), for \( a, b \in \mathbb{R} \) is an \( \mathbb{R} \)-automorphism. In fact, it is an \( \mathbb{R} \)-involution, i.e., an automorphism of order two.

3. If \( F \) is a field, as we always view \( F \subset F[t]/(f) \) for any non-constant polynomial \( f \) in \( F[t] \), \( F[t]/(f) \) is an \( F \)-algebra. [By the definition of \( R \)-algebra \( A \) given in Remark 25.11, we only need a ring homomorphism of a commutative ring \( R \) into the center of a ring \( A \).]

To prove a uniqueness statement for splitting fields, we begin with two lemmas.

**Lemma 49.8.** Let \( \varphi : F \to F' \) be a field isomorphism, \( f \) an irreducible polynomial in \( F[t] \) and \( f' = \tilde{\varphi}(f) \). Then

1. \( f' \) is an irreducible polynomial in \( F'[t] \).
2. There exists a field isomorphism \( \tau : F[t]/(f) \to F'[t]/(f') \) extending \( \varphi \) and taking the root \( \bar{t} = t + (f) \) in \( F[t]/(f) \) to the root \( t = t + (f') \) in \( F'[t]/(f') \).

**Proof.** (1) is clear, so both \( F[t]/(f) \) and \( F'[t]/(f') \) are fields. Let \( \tau : F'[t] \to F'[t]/(f') \) be the canonical epimorphism and \( \mu \) the composition

\[
F[t] \xrightarrow{\tilde{\varphi}} F'[t] \xrightarrow{\tau} F'[t]/(f').
\]

As \( \varphi \) is an isomorphism, so is \( \tilde{\varphi} \). In particular, \( \mu \) is epic. By definition, \( \ker \mu = \{ h \in F[t] \mid \mu(h) = 0 \} = \{ h \in F[t] \mid \tilde{\varphi}(h) \in (f') \} = \{ h \in F[t] \mid h \in (f) \} = (f) \). Therefore, \( \mu \) induces an isomorphism \( \overline{\mu} : F[t]/(f) \to F'[t]/(f') \) given by \( g + (f) \mapsto \tilde{\varphi}(g) + (f') \) by the First Isomorphism Theorem. As \( \overline{\mu} \vert_F = \mu \vert_F = \varphi \) and \( \overline{\mu}(t + (f)) = t + (f') \), the map \( \tau = \overline{\mu} \) works. \( \square \)

**Lemma 49.9.** Let \( K/F \) be an extension of fields and \( f \) an irreducible polynomial in \( F[t] \). If \( \alpha \) in \( K \) is a root of \( f \), then there exists an \( F \)-isomorphism

\[
\theta : F[t]/(f) \to F(\alpha) \text{ given by } t + (f) \mapsto \alpha.
\]

**Proof.** Let \( e_{\alpha} : F[t] \to F[\alpha] \) by \( g \mapsto g(\alpha) \) be the evaluation map at \( \alpha \), a ring epimorphism, and, in fact, an \( F \)-epimorphism. As \( f(\alpha) = 0 \), the element \( \alpha \) is algebraic over \( F \), so, as before, \( e_{\alpha} \) induces an \( F \)-isomorphism \( \theta : F[t]/(m_F(\alpha)) \to F[\alpha] = F(\alpha) \). As \( f \) is irreducible with root \( \alpha \), we have \( (f) = (m_F(\alpha)) \), so \( \theta \) works. \( \square \)
By combining these two lemmas, we shall easily deduce the following:

**Proposition 49.10.** Let $\varphi : F \to F'$ be a field isomorphism, $f$ an irreducible polynomial in $F[t]$ and $f' = \varphi(f) \in F'[t]$. Suppose that $K/F$ is a field extension with $\alpha$ a root of $f$ in $K$ and $K'/F'$ a field extension with $\alpha'$ a root of $f'$ in $K'$. Then there exists an isomorphism $\sigma : F(\alpha) \to F'(\alpha')$ extending $\varphi$ and satisfying $\sigma(\alpha) = \alpha'$.

**Proof.** The two lemmas determine isomorphisms

$$F(\alpha) \xleftarrow{\theta} \frac{F[t]}{(f)} \xrightarrow{\tau} \frac{F'[t]}{(f')} \xrightarrow{\theta'} F'(\alpha')$$

with $\theta(t + (f)) = \alpha$, $\tau(t + (f)) = t + (f')$, and $\theta'(t + (f')) = \alpha'$. The map $\theta' \circ \tau \circ \theta^{-1}$ works. □

The following important corollary follows immediately:

**Corollary 49.11.** Let $K/F$ be an extension of fields, $f$ an irreducible polynomial in $F[t]$ and $\alpha, \beta$ two roots of $f$ in $K$. Then the map

$$\sigma : F(\alpha) \to F(\beta)$$

given by $\sum c_i \alpha^i \mapsto \sum c_i \beta^i$

is an $F$-isomorphism.

The previous proposition is the first step in the induction proof of the result to which we have been striving.

**Theorem 49.12.** Let $\varphi : F \to F'$ be an isomorphism of fields, $f$ a non-constant polynomial in $F[t]$ and $f' = \varphi(f) \in F'[t]$. Suppose that $E$ is a splitting field of $f$ over $F$ and $E'$ is a splitting field of $f'$ over $F'$. Then there exists an isomorphism $\tau : E \to E'$ extending $\varphi$ and taking the set of roots of $f$ in $E$ bijectively onto the set of roots of $f'$ in $E'$ (and preserving multiplicities).

An immediate consequence of this theorem is the result that we had wished to prove.

**Corollary 49.13.** (Uniqueness of Splitting Fields) Let $F$ be a field and $f$ a non-constant polynomial in $F[t]$. If $E$ and $E'$ are splitting fields of $f$ over $F$, then there exists an $F$-isomorphism $E \to E'$ taking the set of roots of $f$ in $E$ bijectively onto the set of roots of $f$ in $E'$ (and preserving multiplicities). In particular, a splitting field of $f$ over $F$ is unique up to an $F$-isomorphism.

Of course, one would like to say given a polynomial there is a unique splitting field. We could do this if we assume that all the fields in which we work lie in some larger field $\Omega$. Then one could say that a field
extension $E/F$ is the splitting field of a non-constant polynomial $f$ in $F[t]$, as the splitting field of $f$ is uniquely determined by the roots of $f$ in $\Omega$. For example, we view all splitting fields of rational polynomials to lie in $\mathbb{C}[t]$.

**Proof.** (of the theorem.) We induct on $n = [E : F]$.

$n = 1$: We have $E = F$, so $f$ splits over $F$, hence $f'$ splits over $F'$. In particular, if $f = a \prod (t - \alpha_i)$ in $F[t]$, then $f' = \varphi(a) \prod (t - \varphi(\alpha_i))$ in $F'[t]$, and the result follows.

$n > 1$: If $f$ does not split over $F$, then $f = gh$, for some polynomials $g$ and $h$ in $F[t]$, with $g$ an irreducible polynomial satisfying $\deg g > 1$. As $\tilde{\varphi} : F[t] \to F'[t]$ is an isomorphism, $f' = g'h'$ with $g' = \tilde{\varphi}(g)$ irreducible of degree greater than one, and $h' = \tilde{\varphi}(h)$. Let $\alpha \in E$ be a root of $g$ in $E$ and $\alpha' \in E'$ be a root of $g'$ in $E'$ (which exist as $f$, $f'$ split over $E$, $E'$, respectively). By Proposition 49.10, there exists an isomorphism $\tau_1 : F(\alpha) \to F'(\alpha')$ extending $\varphi$ and mapping $\alpha$ to $\alpha'$. Let $f = (t - \alpha)f_1$ in $F(\alpha)[t]$. Then $\tilde{\varphi}(f) = (t - \alpha')f_1'$ in $F'(\alpha')[t]$ with $f_1' = \tau_1(f_1)$ and $\tilde{\varphi}(f) = \tilde{\tau}_1(f)$.

We know that $E$ is a splitting field of $f_1$ over $F(\alpha)$ and $E'$ is a splitting field of $f_1'$ over $F'(\alpha')$. Since $[E : F(\alpha)] < [E : F]$, induction provides an isomorphism $\tau : E \to E'$ lifting $\tau_1$, hence $\varphi$, and taking the set of roots of $f_1$ bijectively onto the set of roots of $f_1'$. The result now follows.

The proof above filled in the maps in the following picture:

```
E \xrightarrow{\tau} E' \\
F(\alpha) \xrightarrow{\tau_1} F'(\alpha') \\
F \xrightarrow{\varphi} F'
```

where we shall always use the vertical lines in such pictures as the appropriate set inclusion.

We now look at what we have accomplished, in the following crucial discussion.

**Remark 49.14.** Let $F$ be a field, $f$ a non-constant polynomial in $F[t]$ of degree $n$, and $K$ a splitting field of $f$ over $F$. Let

$$\text{Aut}_F(K) := \{ \tau : K \to K \mid \tau \text{ is an } F\text{-automorphism} \},$$

a group called the *Galois group* of $K$ over $F$. It is also called the *Galois group* of $f$, although if we omit having a fixed splitting field,
then it is only unique up to isomorphism. Now let \( g \) be a non-constant polynomial in \( F[t] \) and set \( S = \{ x \in K \mid g(x) = 0 \} \), i.e., the set of roots of \( g \) in \( K \). Suppose that \( S \) is not empty. Let \( \tau \in \text{Aut}_F(K) \), then \( \tilde{\tau} \) fixes the coefficients of \( g \) in \( F[t] \), so \( \tilde{\tau}(g) = g \). Therefore, if \( x \in K \), we have
\[
\tau(g(x)) = (\tilde{\tau}(g))(\tau(x)) = g(\tau(x)).
\]
Thus, if \( x \in S \), we have
\[
0 = \tau(g(x)) = g(\tau(x)), \text{ i.e., } \tau(x) \in S.
\]
So we have a map,
\[
\text{Aut}_F(K) \to \Sigma(S) \text{ given by the restriction } \tau \mapsto \tau|_S.
\]
Check that this is, in fact, a group homomorphism. Now let \( f = g \) and let \( \alpha_1, \ldots, \alpha_n \) be the roots of \( f \) in \( K \), not necessarily distinct. We know that \( K = F(\alpha_1, \ldots, \alpha_n) = F[\alpha_1, \ldots, \alpha_n] \). Let \( \tau \) and \( \psi \) lie in \( \text{Aut}_F(K) \), and suppose that \( \tau|_S = \psi|_S \). Then \( \tau^{-1} \circ \psi \) fixes \( F \) and \( S \), so \( \tau^{-1} \circ \psi = 1_K \), as \( K = F(S) = F[S] \). We conclude that \( \tau|_S = \psi|_S \) if and only if \( \tau = \psi \). Therefore, the map \( \text{Aut}_F(K) \to \Sigma(S) \) given by \( \tau \mapsto \tau|_S \) is injective, so a group monomorphism, and we can view \( \text{Aut}_F(K) \) as a subgroup of \( \Sigma(S) \) and \( \Sigma(S) \) a subgroup of \( S_n \) with \( \Sigma(S) = S_n \) if and only if \( \alpha_1, \ldots, \alpha_n \) are distinct, i.e., \( f \) has no multiple roots. Suppose now, in addition, that \( f \) is irreducible in \( F[t] \). Fix \( i, j, 1 \leq i, j \leq n \), then \( K \) is a splitting field of \( f \) over \( F(\alpha_i) \) and \( K \) is a splitting field of \( f \) over \( F(\alpha_j) \). We have shown that there exists an \( F \)-isomorphism \( \sigma : F(\alpha_i) \to F(\alpha_j) \) sending \( \alpha_i \) to \( \alpha_j \). The theorem now says that there exists \( \tau \) in \( \text{Aut}_F(K) \) extending \( \sigma \) and taking \( S \) onto \( S \) bijectively. This means that
\[
\text{Aut}_F(K) \subset \Sigma(S) \text{ is a transitive subgroup},
\]
i.e., for every pair of roots \( \alpha_i \) and \( \alpha_j \) in \( S \), there exists a \( \tau \in \text{Aut}_F(K) \) satisfying \( \tau(\alpha_i) = \alpha_j \). We shall see in the sequel, that this will essentially reduce the theory of algebraic extensions to group theory. This is what Galois discovered.

We now give a number of examples to illustrate what we have done.

**Examples 49.15.** Let \( F \) be a field.

1. Let \( f = t^2 + bt + c \) be a polynomial in \( F[t] \) and \( K/F \) an extension field with \( \alpha \) in \( K \) a root of \( f \). If \( \alpha' \) is another root of \( f \) or \( \alpha = \alpha' \) is a multiple root of \( f \), then \( f = (t - \alpha)(t - \alpha') = t^2 - (\alpha + \alpha')t + \alpha\alpha' \) in \( K[t] \). So we have \( b = -\alpha - \alpha' \) and \( c = \alpha\alpha' \). In particular, \( \alpha' = -\alpha - b \) lies in \( K \). If \( \alpha = \alpha' \) is a multiple root, then \( f = (t - \alpha)^2 \) in \( K[t] \) and \( 2\alpha = -b \). If, in addition, \( \text{char} F \neq 2 \), then \( 2 \) is a unit and we have
\[ \alpha = -b/2 \text{ and } f = (t - \alpha)^2 \text{ in } F[t] \text{ splits over } F. \] In the general case, \( f = (t - \alpha)(t - \alpha') \) splits over \( K \) and \( F(\alpha) \) is a splitting field of \( f \) over \( F \), so either \( F = F(\alpha) \) and \( f \) splits over \( F \) or \( [F(\alpha) : F] = 2 \) depending on whether \( f \) is irreducible or not. If \( \text{char } F \neq 2 \), this depends on whether \( b^2 - 4c \) is a square in \( F \) or not, as the quadratic formula can be derived over \( F \) if \( 2 \) is a unit, just as over the real numbers.

2. The field \( \mathbb{Q}(\sqrt{2}, \sqrt{3}) \) is the splitting field of \((t^2 - 2)(t^2 - 3)\) in \( \mathbb{C} \). We already know that \([\mathbb{Q}(\sqrt{2}, \sqrt{3}) : \mathbb{Q}] = 4\) and \(\{1, \sqrt{2}, \sqrt{3}, \sqrt{6}\}\) is a \( \mathbb{Q} \)-basis of \( \mathbb{Q}(\sqrt{2}, \sqrt{3}) \).

3. Let \( f = t^3 - 2 \) in \( \mathbb{Q}[t] \). It is irreducible by Eisenstein’s Criterion. Let \( \omega = (-1 + \sqrt{-3})/2 = \cos(\frac{2}{3}\pi) + \sqrt{-1}\sin(\frac{2}{3}\pi) \) We have the complex conjugate of \( \omega \) satisfies \( \overline{\omega} = \omega^2 \) and \( \omega^3 = 1 \). The roots of \( f \) in \( \mathbb{C} \) are: \( \sqrt[3]{2}, \sqrt[3]{2}\omega, \sqrt[3]{2}\omega^2 \), all distinct. Then \( K = \mathbb{Q}(\sqrt[3]{2}, \sqrt[3]{2}\omega, \sqrt[3]{2}\omega^2) = \mathbb{Q}(\sqrt[3]{2}, \omega) \) is a splitting field of \( f \) over \( \mathbb{Q} \) with \([K : \mathbb{Q}] \leq 3! = 6\). As \( f \) is monic irreducible, we have \( f = m_\mathbb{Q}(\sqrt[3]{2}) \) and, therefore, \( 3 = \deg f = [\mathbb{Q}(\sqrt[3]{2}) : \mathbb{Q}] \mid [K : \mathbb{Q}] \). As \( \omega = \sqrt[3]{2}\omega/\sqrt[3]{2} \) is not real, it does not lie in \( \mathbb{Q}(\sqrt[3]{2}) \). Thus \([K : \mathbb{Q}] > 3\), hence \([K : \mathbb{Q}] = 6\). [Alternatively, \( m_\mathbb{Q}(\omega) = t^2 + t + 1 \) splits in \( K \), so \( \mathbb{Q}(\omega) \subset K \) and \( 2 = \deg m_\mathbb{Q}(\omega) \mid [K : \mathbb{Q}] \), so \([K : \mathbb{Q}] = 6\).] We also have

\[
\begin{array}{c}
K = \mathbb{Q}(\sqrt[3]{2}, \omega) \\
\mathbb{Q}(\sqrt[3]{2}) \quad 2 \quad 3 \\
\mathbb{Q} \quad 6 \quad \mathbb{Q}(\omega) \\
\end{array}
\]

(where the integers in the diagram are the field extension degrees.)

4. Let \( p > 1 \) be a prime and \( f = t^{p-1} + t^{p-2} + \cdots + t + 1 \) in \( \mathbb{Q}[t] \). We know that \( f \) is irreducible by Remark 34.11(3). Let \( \varsigma = e^{2\pi\sqrt{-1}/p} = \cos(2\pi/p) + \sqrt{-1}\sin(2\pi/p) \). Then \( t^p - 1 = (t - 1)f \) is a factorization into irreducibles in \( \mathbb{Q}[t] \). So \( \varsigma^p = 1 \) with \( \varsigma \neq 1 \). As \( f(\varsigma) = 0 \), we have that \([\mathbb{Q}(\varsigma) : \mathbb{Q}] = p - 1 = \varphi(p) \). Similarly, \( \varsigma^2, \ldots, \varsigma^{p-1} \) are roots of \( f \), all distinct and different than \( \varsigma \). It follows that \( \mathbb{Q}(\varsigma) \) is a splitting field of \( f \) (and of \( \varphi(p) \)) over \( \mathbb{Q} \). The polynomial \( f \) is usually denoted by \( \Phi_p \) and called the \( p \)th cyclotomic polynomial in \( \mathbb{Q}[t] \) and satisfies \( \deg \Phi_p = \varphi(p) \).

5. The polynomial \( f = t^4 + t^2 + 1 = (t^2 + t + 1)(t^2 - t + 1) \) in \( \mathbb{Q}[t] \) has roots \( \varsigma = (-1 + \sqrt{-3})/2, \ \zeta = \varsigma^2, \ \eta = (1 + \sqrt{-3})/2, \) and \( \overline{\eta} \),
where $\bar{}$ is complex conjugation. We have $K = \mathbb{Q}(\sqrt{-3}) = \mathbb{Q}(\zeta, \eta)$ is a splitting field of $f$ over $\mathbb{Q}$ and $[K : \mathbb{Q}] = 2$. Note that if $x$ is a root of $f$, then $x^6 = 1$. The polynomial $t^2 - t + 1 = m_\mathbb{Q}(\eta)$ is the 6th cyclotomic polynomial over $\mathbb{Q}$. It is denoted by $\Phi_6$ and satisfies $\deg \Phi_6 = \varphi(6)$.

6. Let $p$ be an odd prime, $r$ a positive integer, and $\zeta = e^{2\pi \sqrt{-1}/p^r}$.

Claim: $\mathbb{Q}(\zeta)$ is the splitting field of $t^{p^r-1}(p-1) + t^{p^r-2}(p-2) + \cdots + t + 1$ and $t^{p^r} - 1 \in \mathbb{Q}[t]$ over $\mathbb{Q}$ and $[\mathbb{Q}(\zeta) : \mathbb{Q}] = \varphi(p^r) = p^r - 1$ (The same is true for $\zeta = e^{2\pi \sqrt{-1}/p^n}$ for any $n$ not divisible by $p$ satisfying $1 \leq n < p^r$). In particular, $\Phi_{p^r} = m_\mathbb{Q}(\zeta)$, the $p^r$th cyclotomic polynomial over $\mathbb{Q}$ satisfies $\deg \Phi_{p^r} = \varphi(p^r)$.

By Exercise 6.14(7), we know that $\varphi(p^r) = p^r - 1(p - 1)$, so it suffices to show $m_\mathbb{Q}(\zeta) = t^{p^r-1}(p-1) + t^{p^r-2}(p-2) + \cdots + t + 1$, as $\zeta^i, 1 \leq i \leq p^r$, are $p^r$ distinct roots of $t^{p^r} - 1$. We need the following identity of binomial coefficients that is easily proven by induction:

$$\sum_{i=j}^{n} \binom{i}{i-j} = \binom{n+1}{n-j}.$$ (49.16)

Let $f = t^{p^r-1}(p-1) + t^{p^r-2}(p-2) + \cdots + t + 1 \in \mathbb{Z}[t]$ and $t = y + 1$.

If we can show that $g(y + 1) = f(t)$ satisfies Eisenstein’s Criterion for $p$, we are done by Remark 34.11(2). By the Children’s Binomial Theorem (Exercise 26.20(7)),

$$(y + 1)^{p^r-1} \equiv (y^{p^r-1} + 1)^i = \sum_{j=0}^{i} \binom{i}{j} y^{p^r-1-j} \mod p\mathbb{Z}[t].$$

We have

$$g(y + 1) = \sum_{i=0}^{p-1} \sum_{j=0}^{i} \binom{i}{j} y^{p^r-1-j} \mod p\mathbb{Z}[t].$$

The coefficient of $y^{p^r-1-j}$ modulo $p\mathbb{Z}[t]$ for $j \neq 0$ is

$$\sum_{i=j}^{p-1} \binom{i}{i-j} = \binom{p}{p-j} \equiv 0 \mod p.$$

The constant term is precisely $p$, so the result follows.

7. Let $f = t^3 + t + 1$ in $(\mathbb{Z}/2\mathbb{Z})[t]$. Check that $f$ has no roots in $(\mathbb{Z}/2\mathbb{Z})$.

Observation 49.17. If $F$ is a field and a non-constant polynomial $g$ in $F[t]$ has no roots in $F$ and is of degree at most three, then $g$ is irreducible in $F[t]$. 
1. If Exercises 49.18.

8. Let $K$ and let $L/K$ such that $f$ has a root $\alpha$ in $L$ and let $K = (\mathbb{Z}/2\mathbb{Z})(\alpha)$. We have $\alpha^3 = -\alpha - 1 = \alpha + 1$ in $K$, so $\alpha^{-1} = \alpha^2 + 1$. Let $f = (t - \alpha)(t^2 + at + b)$ in $K[t]$. Multiplying out, yields $\alpha = a$ and $b = 1 + \alpha^2 = \alpha^{-1}$. It follows that $t^2 + at + b = t^2 + \alpha t + (1 + \alpha^2)$ in $K[t]$. If $e_{\alpha^2} : K[t] \rightarrow K$ is the evaluation map at $\alpha^2$, then

\[
e_{\alpha^2}(t^2 + at + b) = e_{\alpha^2}(t^2 + \alpha t + (1 + \alpha^2)) = \alpha^4 + \alpha^3 + (1 + \alpha^2) = (\alpha^2 + \alpha) + \alpha^3 + (1 + \alpha^2) = (\alpha^2 + \alpha) + (\alpha + 1) + (1 + \alpha^2) = 0.
\]

Consequently, $\alpha^2$ is also a root of $f$ in $K$ with $\alpha \neq \alpha^2$. So we have $f = (t - \alpha)(t - \alpha^2)g$ in $K[t]$ with $\deg g = 1$. Therefore, $f$ splits over $K = (\mathbb{Z}/2\mathbb{Z})(\alpha)$, hence $K$ is a splitting field of $f$ over $(\mathbb{Z}/2\mathbb{Z})$ with $[K : (\mathbb{Z}/2\mathbb{Z})] = \deg f = 3$. Note that $K$ is a field having $2^3 = 8$ elements.

8. Let $p > 0$ be a prime, $F$ a field of characteristic $p$, and $\alpha$ an element in $F$. Then $t^p - \alpha^p = (t - \alpha)^p$ in $F[t]$ by the Children’s Binomial Theorem (Exercise 26.20(7)), has $F$ as a splitting field of $f$ and $f$ has only one distinct root of multiplicity $p$. Now let $L/F$ be an extension field, $\alpha$ an element in $L$ satisfying $a = \alpha^p$ lies in $F$. Then $K = F(\alpha)$ is a splitting field of $t^p - a$ in $F[t]$ over $F$.

Claim. If $\alpha$ is not an element in $F$, then $[K : F] = p$.

Suppose that $t^p - \alpha = fg$ in $F[t]$ with $f$ and $g$ lying in $F[t]$. We may assume that both $f$ and $g$ are monic and $f$ not a unit. As $fg = t^p - a = (t - \alpha)^p$ in the UFD $K[t]$, we must have $f = (t - \alpha)^r$ in $K[t]$ for some $r$ satisfying $1 \leq r \leq p$. As $f$ lies in $F[t]$, its constant term $\pm \alpha^r$ must lie in $F$. If $r < p$, then $r$ and $p$ are relatively prime, so $\alpha^r$, $\alpha^p$ lying in $F$ implies that $\alpha$ lies in $F$, a contradiction. Therefore, $g = 1$ and $r = p$, so $t^p - a$ is irreducible in $F[t]$.

Exercises 49.18.

1. If $f \in \mathbb{Q}[t]$ and $K$ is a splitting field of $f$ over $\mathbb{Q}$, determine $[K : \mathbb{Q}]$ if $f$ is:
   (a) $t^4 + 1$.
   (b) $t^6 + 1$.
   (c) $t^4 - 2$.
   (d) $t^6 - 2$.
   (e) $t^6 + t^3 + 1$.

2. Find the splitting fields $K$ for $f \in \mathbb{Q}[t]$ and $[K : \mathbb{Q}]$ if $f$ is:
   (a) $t^4 - 5t^2 + 6$.
(b) \( t^6 - 1 \).
(c) \( t^6 - 8 \).

3. Show both of the following:
   (a) If \( K / \mathbb{Q} \) and \( \sigma \in \text{Aut} K \) then \( \sigma \) fixes \( \mathbb{Q} \).
   (b) The fields \( \mathbb{Q}(\sqrt{2}) \) and \( \mathbb{Q}(\sqrt{3}) \) are not isomorphic.

4. Prove Identity 49.16.

5. A primitive \( n \)th root of unity is an element \( z \in \mathbb{C} \) such that \( z^n = 1 \) and \( z^r \neq 1 \) for \( 1 \leq r < n \). Show the following:
   (a) There exist \( \varphi(n) := |\{d \mid 0 \leq d \leq n, (d,n) = 1\}| \) primitive \( n \)th roots of unity.
   (b) If \( \omega \) is a primitive \( n \)th root of unity, then \( \mathbb{Q}(\omega) \) is a splitting field of \( t^n - 1 \in \mathbb{Q}[t] \).
   (c) Let \( \omega_1, ..., \omega_{\varphi(n)} \) be the \( \varphi(n) \) primitive \( n \)th roots of unity of \( t^n - 1 \in \mathbb{Q}[t] \) and \( \sigma \in \text{Aut} \mathbb{Q}(\omega_1) \) then \( \sigma(\omega_1) = \omega_i \) for some \( i, 1 \leq i \leq \varphi(n) \).

6. Continued from Exercise 5. Show:
   (a) Let \( \Phi_n(t) = (t - \omega_1) \cdots (t - \omega_{\varphi(n)}) \). Then show \( \Phi_n(t) \in \mathbb{Q}[t] \).
   (b) \( \Phi_n(t) \in \mathbb{Z}[t] \).

7. Continued from Exercise 6. Show:
   (a) \( \Phi_n(t) \in \mathbb{Z}[t] \) is irreducible.
   (b) Calculate \( \Phi_n(t) \) for \( n = 3, 4, 6, 8 \) explicitly and show directly that \( \Phi_n(t) \in \mathbb{Z}[t] \) is irreducible.

8. Let \( F = \mathbb{Z}/p\mathbb{Z} \). Show all of the following:
   (a) There exists a polynomial \( f \in F[t] \) satisfying \( \deg f = 2 \) and \( f \) irreducible.
   (b) Use the \( f \) in (a) to construct a field with \( p^2 \) elements.
   (c) If \( f_1, f_2 \in F[t] \) have \( \deg f_i = 2 \) and \( f_i \) irreducible for \( i = 1, 2 \), show that their splitting fields are isomorphic.

9. Let \( F \) be a field of characteristic different from two and \( f = at^2 + bt + c \) in \( F[t] \). State and prove the quadratic formula.

10. Let \( K/F \) be an extension of fields and \( f \in F[t] \).
    (a) If \( \varphi : K \to K \) is an \( F \)-automorphism, then \( \varphi \) takes roots of \( f \) in \( K \) to roots of \( f \) in \( K \).
    (b) If \( F \subset \mathbb{R} \) and \( \alpha = a + b\sqrt{-1} \) is a root of \( f \) with \( a, b \in \mathbb{R} \) then \( \overline{\alpha} = a - b\sqrt{-1} \) is also a root of \( f \).
    (c) Let \( F = \mathbb{Q} \). If \( m \in \mathbb{Z} \) is not a square and \( a + b\sqrt{m} \in \mathbb{C} \) is a root of \( f \) with \( a, b \in \mathbb{Q} \) then \( a - b\sqrt{m} \) is also a root of \( f \) in \( \mathbb{C} \).

11. Any (field) automorphism \( \varphi : \mathbb{R} \to \mathbb{R} \) is the identity automorphism.
12. Let \( f = (t^2 - p_1) \cdots (t^2 - p_n) \) in \( \mathbb{Q}[t] \) with \( p_1, \ldots, p_n \) be \( n \) distinct prime numbers. Show that \( K = \mathbb{Q}[\sqrt{p_1}, \ldots, \sqrt{p_n}] \) is a splitting field of \( f \) over \( \mathbb{Q} \) and \([K : \mathbb{Q}] = 2^n\). Formulate a generalization of the statement for which your proof still works.

13. Find a splitting field of \( f \in F[t] \) if \( F = \mathbb{Z}/p\mathbb{Z} \) and \( f = t^e - t, \ e > 0 \).

14. Let \( K \) be a splitting field over \( F \) of \( f \) in \( F[t] \setminus F \) and \( g \) an irreducible polynomial in \( F[t] \). Suppose that \( g \) has a root in \( K \). Show that \( g \) splits over \( K \).

15. Suppose \( F \) is a field and \( K \) is a splitting field of \( f \in F[t] \) over \( F \). If \( K/E/F \) is an intermediate field, \( \sigma : E \to K \) be a monomorphism fixing \( F \), show that there exists an \( F \)-automorphism \( \hat{\sigma} : K \to K \) lifting \( \sigma \), i.e., such that \( \hat{\sigma}|_E = \sigma \).

16. Let \( F \) be a field of characteristic \( p > 0 \). Show that \( f = t^4 + 1 \in F[t] \) is not irreducible. Let \( K \) be a splitting field of \( f \) over \( F \). Determine which finite fields \( F \) satisfies \( K = F \).

17. Let \( f = t^6 - 3 \in F[t] \). Construct a splitting field \( K \) of \( f \) over \( F \) and determine \([K : F]\) for each of the cases: \( F = \mathbb{Q}, \mathbb{Z}/5\mathbb{Z}, \) or \( \mathbb{Z}/7\mathbb{Z} \). Do the same thing if \( f \) is replaced by \( g = t^6 + 3 \in F[t] \) for the same fields \( F \).

### 50. Algebraically Closed Fields

Recall that a field \( K \) is called \emph{algebraically closed} if every non-constant polynomial in \( K[t] \) has a root in \( K \), equivalently that every non-constant polynomial in \( K[t] \) splits over \( F \). Clearly, this is also equivalent to the condition that if \( L/K \) is an algebraic extension of fields, then \( L = K \). We call an extension field \( K \) of \( F \) an \emph{algebraic closure of \( F \) if \( K/F \) is algebraic and \( K \) is algebraically closed. We shall later prove the Fundamental Theorem of Algebra that \( \mathbb{C} \) is an algebraically closed field. We know that \( \mathbb{C} \) is not an algebraic closure of \( \mathbb{Q} \). In fact, assuming the Fundamental Theorem of Algebra, we showed in Example 47.22(4) that \( \Omega := \{ x \in \mathbb{C} \mid x \text{ algebraic over } \mathbb{Q} \} \) is an algebraic closure of \( \mathbb{Q} \) and is countable. In this section, we shall show that given any field \( F \), there exists an algebraic closure of \( F \), and it is unique up to an \( F \)-homomorphism. By the argument in Example 47.22(5), to establish the existence of an algebraic closure of \( F \), it suffices to show that \( F \) is a subfield of some algebraically closed field.

We know that if we have a finite set of polynomials, \( \{f_1, \ldots, f_n\} \), then we can construct a field \( K \) such that each \( f_i \) splits in \( K \) and \( K \) is the smallest such field by letting \( K \) be the splitting field of \( \prod f_i \) over \( F \). Another way of constructing this \( K \) would be to find a splitting field
of $f_1$ over $F$, then $K_2$ a splitting field of $f_2$ over $K_1$, and continue in this way. We use both for the case of an infinite set of polynomials.

To use these two ideas, we first set up some notation. Let $I$ be an indexing set and $S = \{t_i \mid i \in I\}$ be a set of distinct indeterminants. Then $F[S]$ is the set of polynomials, coefficients in $F$ in the $t_i$'s, $i \in I$. If $J \subseteq I$, let $S_J = \{t_j \mid j \in J\}$. As each such polynomial needs only finitely many indeterminants to determine it, i.e., $f \in F[S_J]$, for some finite subset $J \subseteq I$, we have

$$F[S] = \bigcup_{J \subseteq I, \text{finite}} F[S_J].$$

Note that $F[S]$ is a domain, and in fact, a UFD. Using this notation, we give Artin’s proof of the following theorem.

**Theorem 50.1.** Let $F$ be a field. Then there exists an algebraically closed field $K$ containing $F$.

**Proof.** Let $I = F[t] \setminus F$ be our indexing set and $S = \{t_f \mid f \in I\}$. So if $f \in F[t]$, we have $f(t_f)$ lies in $F[t_f] \subset F[S]$. Set $T = \{f(t_f) \mid f \in F[t] \setminus F\}$ and $\mathfrak{A}$ the ideal in $F[S]$ generated by the set $T$. So

$$\mathfrak{A} = \left\{ \sum_{\text{finite}} g_i f_i(t_f) \mid g_i \in F[S] \right\}.$$

**Claim.** $\mathfrak{A} \subset F[S]$:

Suppose that $\mathfrak{A} = F[S]$. Then there exist $f_i \in T$ and $g_i \in F[S]$, $i = 1, \ldots, n$, some $n$, satisfying

$$1 = g_1 f_1(t_{f_1}) + \cdots + g_n f_n(t_{f_n})$$

and a finite subset $J \subseteq I$ with $g_1, \ldots, g_n$ lying in $F[S_J]$. Relabel each $t_{f_i}$ as $t_i$ and the $t_j$'s in $J$ such that $J = \{t_1, \ldots, t_N\}$ some $N$, so $f_i, g_i \in F[t_1, \ldots, t_N]$. We know that there exists a field extension $K/F$ such that each $f_i$ has a root $\alpha_i$ in $K$ for $1 \leq i \leq n$ by the preliminary remarks. Set $\alpha_i = 0$ for $n < i \leq N$. Applying the evaluation map $e_{\alpha_1, \ldots, \alpha_n}$ yields

$$0 = \sum_{\text{finite}} g_i(\alpha_1, \ldots, \alpha_n) f_i(\alpha_i) = 1$$

in $K$, a contradiction. This establishes the claim.

By the claim, there exists a maximal ideal $\mathfrak{m}$ satisfying $\mathfrak{A} \subseteq \mathfrak{m} < F[S]$ (which we know exists by Zorn’s Lemma). Define $L_1 := F[S]/\mathfrak{m}$, viewed as a field extension of $F$ in the usual way (i.e., identifying $F$ and $\varphi(F)$ where $\varphi : F[S] \to L$ is given by $g \mapsto g + \mathfrak{m}$). By construction, if $f$ is a non-constant polynomial in $F[t]$, then $f$ has a root $t_f + \mathfrak{m}$ in $L_1$. Inductively, define $L_i$ such that $L_i/L_{i-1}$ and every non-constant
polynomial \( g \in L_{i-1}[t] \) has a root in \( L_i \) by repeating the same construction. Set \( L = \cup L_i \). If \( \alpha, \beta \) lie in \( L \) with \( \beta \) nonzero, then \( \alpha \pm \beta, \alpha \beta, \beta^{-1} \) lie in some subfield \( L_i \) of \( L_i \) hence \( L \) is a field. Finally, let \( f \) be a non-constant polynomial in \( L[t] \). As \( f \) has finitely many nonzero coefficients, there exists an \( i \) such that \( f \) lies in \( L_i[t] \) hence has a root in \( L_{i+1} \subset L \). Therefore, \( L \) is algebraically closed.

We showed in the argument in Example 47.22(5) that if \( F \) is a field and \( K/F \) a field extension with \( K \) algebraically closed, then the elements in \( K \) algebraic over \( F \) form an algebraic closure of \( F \). Therefore, we have established the following result.

**Corollary 50.2.** Let \( F \) be a field. Then an algebraic closure of \( F \) exists.

Algebraic closures are “super” splitting fields, so we expect a uniqueness statement. We prove this.

**Theorem 50.3.** Let \( K/F \) be an algebraic extension of fields and \( L \) an algebraically closed field. Suppose that there exists a homomorphism \( \sigma : F \to L \). Then there exists a homomorphism \( \tau : K \to L \) that lifts \( \sigma \). In particular, if \( K \) is also algebraically closed and \( L/\sigma(F) \) is algebraic, then \( \tau \) is an isomorphism.

**Proof.** We use a Zorn Lemma argument that is often repeated in algebra to define maps. Let

\[
S := \{(E, \eta) \mid K/E/F \text{ an intermediate field} \quad \eta : E \to L \text{ a homomorphism lifting } \sigma \}.
\]

Partially order the set \( S \) by \( \leq \) defined by

\[
(E, \eta) \leq (E', \eta') \text{ if } E \subseteq E' \text{ and } \eta'|_E = \eta.
\]

As \( (F, \sigma) \) lies in \( S \), the set \( S \) is nonempty. Let \( C = \{(E_i, \eta_i) \mid i \in I\} \) be a chain in \( S \). We know that \( E = \cup E_i \) is a field. Define

\[
\eta : E \to L
\]

as follows: If \( \alpha \) lies in \( E \), choose \( i \in I \) such that \( \alpha \in E_i \) and set \( \eta(\alpha) := \eta_i(\alpha) \). By the definition of our partial ordering, \( \eta(\alpha) \) is independent of \( i \), i.e., \( \eta \) is well-defined. It follows that \( (E, \eta) \) is an upper bound for \( C \). By Zorn’s Lemma, there exists a maximal element \( (E, \eta) \) in \( S \), so \( K/E/F \) and \( \eta : E \to L \) extends \( \sigma \).

**Claim.** \( K = E \):

Suppose this is false and \( \alpha \in K \setminus E \). As \( \alpha \) is algebraic over \( F \), it is algebraic over \( E \). Let \( \beta \) be a root of \( \tilde{\eta}(m_E(\alpha)) \) in \( \eta(E)[t] \) in the
algebraically closed field \(L\). Then we know by Proposition 49.10 that there exists an isomorphism
\[
\eta' : E(\alpha) \to \eta(E)(\beta) \quad \text{with} \quad \alpha \mapsto \beta \quad \text{lifting} \quad \eta.
\]
It follows that \((E(\alpha), \eta')\) lies in \(S\) with \((E, \eta) < (E(\alpha), \eta')\), contradicting the maximality of \((E, \eta)\). This proves the claim.

If \(K\) is algebraically closed so is the subfield \(\tau(K)\) in \(L\). If \(L/\sigma(F)\) is algebraic, then \(L/\tau(K)\) is algebraic, hence \(L = \tau(K)\) and the theorem if proven. \(\square\)

We immediately see that this implies the desired result.

**Corollary 50.4.** Let \(F\) be a field. Then any algebraic closure of \(F\) is unique up to an \(F\)-isomorphism.

Because of the corollary, one usually fixes an algebraic closure \(K\) of a field \(F\) and views all roots of all polynomials in \(F[t]\) as lying in \(K\). More generally, one fixes some \(E/F\) with \(E\) algebraic closed and only works with fields \(K\) with \(E/K/F\). Then \(E\) contains a unique algebraic closure of \(F\), hence for each polynomial in \(F[t]\) a unique splitting field in \(E\). One then views every algebraic extension of \(F\) to be in \(E\). For example, one can let \(E = \mathbb{C}\) when \(F = \mathbb{Q}\). More generally, if \(F\) is a field of characteristic zero and \(\tilde{F}\) an algebraically closed field, we can view \(F/\mathbb{Q}\) and the algebraic closure of \(\mathbb{Q}\) in \(\tilde{F}\) in \(\mathbb{C}\) by replacing it with an isomorphic copy. We shall always do so implicitly.

**Exercises 50.5.**

1. Let \(F\) be a finite field and \(K\) an algebraic closure of \(F\). Show that \(K\) is countable and infinite. (In all other cases, an algebraic closure of a field has the same cardinality of the given field).

2. Let \(S\) be a collection of polynomials over a field \(F\). Define what it means for an extension field \(K/F\) to be a splitting field of the collection \(S\) and prove that one exists and is unique up to an \(F\)-isomorphism. Also prove that an algebraic closure of \(F\) is a splitting field for the collection of irreducible polynomials in \(F[t]\).

**51. Constructible Numbers**

The ancient Greeks were very interested in constructing geometric objects under certain rules. The most famous ones that they could not solve were the following:

**Euclidean Construction Problems.** Using a straight-edge and compass, can you:
51. CONSTRUCTIBLE NUMBERS

I. Trisect a given angle?

II. (Delian Problem) Double a cube, i.e., construct the edge of a cube whose volume is twice the volume of a given cube?

III. Square a circle, i.e., construct a square whose area equals the area of a given circle?

IV. Construct a regular \( n \)-gon, \( n \geq 3 \)?

Given a specific set of initial points, the following constructions were allowed:

\begin{itemize}
  \item C1. You can draw a line through two initial points.
  \item C2. You can draw a circle with center at an initial point with radius equal to the length between two initial points.
\end{itemize}

From these constructions on your given set of initial points, you were allowed to enlarge this original set by adding the points of intersection arising from the constructions C1 and C2 and add further points by iterating this process.

We shall solve the Euclidean Problems, by converting them into problems in field theory. In the case of squaring the circle, some analysis will also be needed.

**Algebraic Reformulation** (Preliminaries). Let \( S = \{P_1, \ldots, P_n\} \) be a finite non-empty set of points in the euclidean plane. Recursively define \( S_r \) as follows:

\( S_0 = S \) and having defined \( S_r \) define \( S_{r+1} := S_r \cup T_r \)

where \( T_r \) is the union of the following sets:

1. The set of points of intersections of lines through two distinct points in \( S_r \) (that we call \( S_r \)-lines).
2. The set of points of the intersection of circles having center in \( S_r \) and radii equal in length to line segments connecting two distinct points in \( S_r \) (that we call \( S_r \)-circles).
3. The set of intersections of \( S_r \)-circles and \( S_r \)-lines.

Now set

\[
C(P_1, \ldots, P_n) := \bigcup_{r=0}^{\infty} S_r,
\]

the set of points constructible from \( P_1, \ldots, P_n \). We say a point \( P \) in the euclidean plane is constructible from \( P_1, \ldots, P_n \), if \( P \in C(P_1, \ldots, P_n) \) and not constructible otherwise. If \( C = C(P_1, \ldots, P_n) \), call a line segment a \( C \)-line if an \( S_r \)-line for some \( r \) and a \( C \)-circle if an \( S_r \)-circle for some \( r \).
Remarks 51.1. Let \( P_1, \ldots, P_n \) be our initial given points in the euclidean plane and \( C = C(P_1, \ldots, P_n) \) Using elementary euclidean geometry, we deduce the following:

1. \( C(P_1) = \{ P_1 \} \). As this is not interesting, we shall always assume that we begin with at least two given points, i.e., \( n \geq 2 \).

2. We can begin to define a coordinate system of our euclidean plane by setting \( P_1 = (0, 0) \), the origin, and \( P_2 = (1, 0) \).

3. Given a line segment of length \( d \) between two points in \( C \), we can construct a line segment of length \( md \) for all integers \( m \), can bisect the given line segment between two points in \( C \), and can construct a perpendicular at each endpoint of the line segment between two points in \( C \). In particular, we can construct axes.

4. We can drop a perpendicular from a point in \( C \) to a given \( C \)-line as well as construct a parallel through a point in \( C \) to a \( C \)-line. In particular, we see that any point in \( \mathbb{Z}^2 \) is constructible from \( P_1, \ldots, P_n \).

5. Given two \( C \)-lines of angles \( \theta \) and \( \varphi \) with the \( X \)-axis respectively, we can construct rays (\( C \)-lines through the origin) of angles \( \theta \pm \varphi \), \(-\theta \), \( \theta/2 \). In fact, this can be done with the \( X \)-axis replaced by any allowable \( C \)-line.

It is now convenient to complexify everything, i.e., to replace the euclidean plane with the complex plane. Let \( P_1, \ldots, P_n \), \( n \geq 2 \), be points in the euclidean plane. As above, we let \( P_1 = (0, 0) \) and \( P_2 = (1, 0) \). To complexify, if \( P_i = (x_i, y_i) \) lie in \( \mathbb{R}^2 \), let \( z_i = x_i + y_i \sqrt{-1} \), e.g., we have \( z_1 = 0 \) and \( z_2 = 1 \), and let \( C(z_1, \ldots, z_n) \) denote the complex numbers constructible by straight-edge and compass from \( z_1, \ldots, z_n \).

We call \( C(z_1, \ldots, z_n) \) the set of complex constructible numbers from \( z_1, \ldots, z_n \). We have \( (x, y) \in C(P_1, \ldots, P_n) \) if and only if \( x + y \sqrt{-1} \in C(z_1, \ldots, z_n) \). Note that we already know that \( C(z_1, \ldots, z_n) \) contains the Gaussian integers. But we have much more.

Theorem 51.2. Let \( n \geq 2 \) and \( z_1 (= 0), z_2 (= 1), \ldots, z_n \) be complex numbers. Then \( C(z_1, \ldots, z_n) \) is a field satisfying \( \mathbb{Q} \subset C(z_1, \ldots, z_n) \subset \mathbb{C} \).

Proof. As \( C(z_1, \ldots, z_n) \subset \mathbb{C} \) with \( \text{char } \mathbb{C} = 0 \), it suffices to show that \( C(z_1, \ldots, z_n) \) is a field. Let \( z, z' \) lie in \( C(z_1, \ldots, z_n) \), with \( z \) nonzero. As we can draw parallel lines, we can add given vectors (line segments starting at the origin), so we see that \( z \pm z' \) lies in \( C(z_1, \ldots, z_n) \). We need to show that \( zz' \) and \( z^{-1} \) lie in \( C(z_1, \ldots, z_n) \). Write \( z \) and \( z' \) in polar form, say \( z = re^{\sqrt{-1} \theta} \) and \( z' = r'e^{\sqrt{-1} \theta'} \) with \( r, r', \theta, \theta' \) real numbers, \( r, r' \) non-negative. Then \( zz' = rr' e^{\sqrt{-1} (\theta + \theta')} \) and \( z^{-1} = r^{-1} e^{-\sqrt{-1} \theta} \).
We know that we can construct rays of angle \( \theta \pm \theta' \) and \(-\theta\) and line segments of length \( r = |z| \) and \( r' = |z'| \). On the \( X\)-axis in the complex plane mark off 1 and \( 1 + r \) and the ray of angle \( \theta' \) through \( z' \), which we may assume to be nonzero. Then the line through \( 1 + r \) parallel to the line through \( z' \) and 1 intersects the ray of an angle \( \theta' \) with the \( X\)-axis at \( (r + x)e^{\sqrt{-1}\theta'} \), and \( x \) is seen to be \( rr' \) by similar triangles. Reversing, the roles of \( x \) and \( r' \) constructs \( r/r' \). (Cf. Figure 51.3 below.) As associativity, commutativity, etc. hold in \( \mathbb{C} \), we see that \( C(z_1, \ldots, z_n) \) is a field. \( \square \)

The field \( C(z_1, \ldots, z_n) \) is easily characterized using computation of the intersection points of lines and circles.

**Theorem 51.4.** Let \( z_1 (= 0), z_2 (= 1), \ldots, z_n \) be points in \( \mathbb{C} \) with \( n \geq 2 \). Then the field \( C(z_1, \ldots, z_n) \) has the following properties:

1. The point \( z_i \) lies in \( C(z_1, \ldots, z_n) \) for \( i = 1, \ldots, n \).
2. If \( z \) lies in \( C(z_1, \ldots, z_n) \), so does its complex conjugate \( \overline{z} \).
3. If \( z^2 \) lies in \( C(z_1, \ldots, z_n) \), so does \( z \), i.e., if \( z \) lies in \( C(z_1, \ldots, z_n) \), so does \( \sqrt{z} \).

Moreover, \( C(z_1, \ldots, z_n) \) is the unique smallest subfield of \( \mathbb{C} \) satisfying these properties, i.e., if a subfield \( C' \) of \( \mathbb{C} \) satisfies (1), (2), and (3), then \( C(z_1, \ldots, z_n) \subset C' \).

**Proof.** (1) is clear.

(2): If \( z = re^{\sqrt{-1}\theta} \), then \( \overline{z} = re^{-\sqrt{-1}\theta} \), and (2) follows.

(3): Since we can construct rays of angle \( \theta/2 \) from rays of angle \( \theta \), it suffices to show that we can construct \( \sqrt{r} \) given a segment of length \( r \). Draw a circle centered at the point \( (1 + r)/2 \) on the \( X\)-axis of radius \( (1 + r)/2 \). Erect a perpendicular at the point 1 on the \( X\)-axis parallel to the \( Y\)-axis. This line segment has length \( \sqrt{r} \) by the Pythagorean Theorem. (Cf. Figure 51.5.)
This shows that $C(z_1, \ldots, z_n)$ satisfies the desired properties. Suppose that $C'$ also does. As $z_1, \ldots, z_n$ all lie in $C'$, by the definition of $C(z_1, \ldots, z_n)$, it suffices to show that intersection points of two $C'$-lines, two $C'$-circles, and a $C'$-line and a $C'$-circle, lie in $C'$. To do this note that $\sqrt{-1}$ lies in $C'$ by (3) and if $z = a + b\sqrt{-1}$ with $a, b$ in $\mathbb{R}$, then $z$ lies in $C'$ if and only if $a$ and $b$ lie in $C' \cap \mathbb{R}$, e.g., if $z$ lies in $C'$, then by (2), so does $b = -\sqrt{-1}(z - \bar{z})/2$ as $C'$ is a field. This means that the equations defining $C'$-lines and $C'$-circles have coefficients in $C' \cap \mathbb{R}$, so solving these equations leads to solutions in $C'$ (using (3) for intersections with circles).

We wish to further characterize the field theoretic properties of fields of constructible numbers. Property (3) in the theorem gives the key. Such fields are closed under taking square roots of elements in the field.

**Definition 51.6.** Let $K/F$ be an extension of fields. We call $K/F$ a *square root tower over $F$* if there exist elements $u_1, \ldots, u_n$ in $K$ for some positive integer $n$, satisfying:

1. $K = F(u_1, \ldots, u_n)$.
2. $u_1^2$ lies in $F$.
3. $u_i^2$ lies in $F(u_1, \ldots, u_{i-1})$ for each $i > 1$.

**Remarks 51.7.** Let $K/F$ be a square root tower defined by $K = F(u_1, \ldots, u_n)$ satisfying the above properties and $E/F$ an extension field. Assume that both $K$ and $E$ lie in some larger field.

1. We have

$$[F(u_1, \ldots, u_i) : F(u_1, \ldots, u_{i-1})] = \begin{cases} 1 & \text{if } u_i \in F(u_1, \ldots, u_{i-1}) \\ 2 & \text{otherwise}. \end{cases}$$

In particular, $[F(u_1, \ldots, u_i) : F] = 2^e$ for some $e \leq n.$
2. The extension $E(K)/E$ is a square root tower. (Of course, $E(K) = E(u_1, \ldots, u_n)$.)
3. If $E/F$ is also a square root tower, then so is $E(K)/F$.
4. If $E/K$ is a square root tower, then so is $E/F$.
5. Suppose that $K \subset \mathbb{C}$ and $\overline{F} := \{ \overline{z} \mid z \in F \}$, where $\overline{z}$ is the complex conjugate of $z$. Then $\overline{F}$ is a field and $\overline{K}/\overline{F}$ is a square root tower, with $\overline{K} = \overline{F}(\overline{z_1}, \ldots, \overline{z_n})$.
6. Suppose that $\text{char } F \neq 2$ and $E/F$ satisfies $[E:F] \leq 2$. Then $E/F$ is a square root tower:
   
   We may assume that $[E:F] = 2$, so there exists an $\alpha \in E \setminus F$ satisfying $E = F(\alpha)$. Let $m_F(\alpha) = t^2 + bt + a$ in $F[t]$. Then $m_F(\alpha)$ splits over $E$ and has roots $\alpha$ and $\alpha_1 = -b - \alpha$. Let $\beta = 2\alpha + b$. As $\alpha^2 + b\alpha = -a$ and $\beta$ does not lie in $F$, we see that $E = F(\beta)$ and $\beta^2 = -4a + b^2$ lies in $F$.

   In general, it is not true if a field extension $E/F$ has degree a power of two, then it is a square root tower. [Can you find an example?] However, we shall prove the following theorem later (cf. Theorem 56.11):

**Theorem 51.8.** (Square Root Tower Theorem) Let $F$ be a field of characteristic zero and $K/F$ an extension of fields that is a splitting field of some non-constant polynomial in $F[t]$. Then $K/F$ is a square root tower if and only if $K/F$ is a field extension of degree a power of two.

   In fact, we shall prove a stronger result, that will be applicable to fields of arbitrary characteristic. We shall use this result when needed later in this section.

   The notion of square root tower allows us to give a sufficient condition for a complex number to be constructible from a set of complex numbers.

**Theorem 51.9.** (Constructibility Criterion) Given complex numbers $z_1 (= 0), z_2 (= 1), \ldots, z_n$ with $n \geq 2$, set $F = \mathbb{Q}(z_1, z_2, \ldots, z_n, \overline{z_1}, \ldots, \overline{z_n})$. If $z$ is a complex number, then $z$ is constructible from $z_1, z_2, \ldots, z_n$, i.e., $z$ lies in $C(z_1, z_2, \ldots, z_n)$ if and only if there exists a square root tower $K/F$ (with $K \subset \mathbb{C}$) satisfying $z \in K$.

   **Proof.** Let $C' = \{ z \in \mathbb{C} \mid \text{There exists a square root tower } K/F \text{ with } z \in K \}$.

   By definition, $F \subset C'$. We must show that $C' \subset C(z_1, \ldots, z_n)$. We know that $C(z_1, \ldots, z_n)$ is a field containing $z_1, \ldots, z_n$ closed under complex conjugation and extracting square roots. In particular, $F$ is a
subfield of $C(z_1, \ldots, z_n)$ as is any square root tower $K/F$. In particular, $C' \subset C(z_1, \ldots, z_n)$. To show that $C(z_1, \ldots, z_n) \subset C'$, it suffices to show by Theorem 51.4 that $C'$ is a field closed under complex conjugation and extracting square roots. Let $z$ and $z'$ be elements in $C'$ satisfying $z \in K$ and $z' \in K'$ with $K/F$ and $K'/F$ square root towers. Then $K(K')/F$ is a square root tower by Remark (3) above containing $zz'$, $z \pm z'$, and $z^{-1}$ (if $z \neq 0$). As $K(K') \subset C'$, we have $C'$ is a field. Moreover, $K(\sqrt{z})/K$ is a square root tower, hence $K(\sqrt{z})/F$ is also a square root tower as $K/F$ is. Since $F = \mathbb{F}$, we also have $\mathbb{F}/F$ is a square root tower. Therefore, $\sqrt{z}$ and $z$ lie in $C'$ and $C(z_1, \ldots, z_n) \subset C'$ as needed.

**Corollary 51.10.** Let $z_1 (= 0), z_2 (= 1), \ldots, z_n$ be complex numbers with $n \geq 2$ and $z$ a complex number in $C(z_1, \ldots, z_n)$. Then $z$ is algebraic over $F = \mathbb{Q}(z_1, z_2, \ldots, z_n, \sqrt{z_1}, \ldots, \sqrt{z_n})$ and $[F(z) : F]$ is a power of two.

**Proof.** As $z$ lies in $C(z_1, \ldots, z_n)$, there exists a square root tower $K/F$ with $z \in K$, hence $[F(z) : F] \mid [K : F]$, a power of two.

In many euclidean constructions, one starts with two points, $P_1$ and $P_2$, which we may assume correspond to complex numbers $z_1 = 0$ and $z_2 = 1$. In this case, $\mathbb{Q} = \mathbb{Q}(z_1, z_2, \overline{z_1}, \overline{z_2})$ and $C(z_1, z_2)$ is called the field of (complex euclidean) constructible numbers and any $z \in C(z_1, z_2)$ is called constructible. So if $z$ is a constructible number, we have $[\mathbb{Q}(z) : \mathbb{Q}]$ is a power of two.

We now can solve the (Greek) Euclidean Construction Problems. Recall that a cubic (or quadratic) polynomial in $F[t]$, $F$ a field, is irreducible if and only if it does not have a root in $F$.

**Trisection of an angle $\theta$:** To solve this construction problem, we need the following lemma.

**Lemma 51.11.** Let $\alpha$ be a ray with the $X$-axis as one side. Then it is constructible form $z_1 (= 0), z_2 (= 1), \ldots, z_n$ if and only if $e^{\sqrt{-1} \alpha}$ is constructible from $z_1 (= 0), z_2 (= 1), \ldots, z_n$ if and only if $\cos \alpha$ is constructible from $z_1 (= 0), z_2 (= 1), \ldots, z_n$.

**Proof.** As $z = e^{\sqrt{-1} \alpha} = \cos \alpha + \sqrt{-1} \sin \alpha$ and we can erect and drop perpendiculars, this is immediate.

The trisection of an angle $\theta$ is, therefore, equivalent to:

**Problem.** Given $z_1 = 0, z_2 = 1, z_3 = \cos \theta$, can we construct $\cos(\theta/3)$ from $z_1 = 0, z_2 = 1, z_3$?
Let \( F = \mathbb{Q}(\cos \theta, \cos \theta) = \mathbb{Q}(\cos \theta) \). For any angle \( \alpha \), we have the trigonometric identity
\[
\cos 3\alpha = 4 \cos^3 \alpha - 3 \cos \alpha.
\]
So \( \beta = \cos(\theta/3) \) is a root in \( \mathbb{C} \) of \( f = 4t^3 - 3t - \cos \theta \) in \( F[t] \). The polynomial is reducible, i.e., not irreducible in \( F[t] \) if and only if \( f \) has a root in \( F \) if and only if for any root \( \gamma \) of \( f \) in \( \mathbb{C} \), we have \( \deg m_F(\gamma) < 3 \) if and only if \( F(\beta)/F \) is a square tower if and only if \( \beta \) is constructible from \( z_1, z_2, z_3 \). Hence \( \beta \not\in \mathbb{C}(z_1, z_2, z_3) \) if and only if \( f \) has no root in \( F \). We, therefore, have proven the following:

**Theorem 51.12.** An angle \( \theta \) can be trisected with straight-edge and compass if and only if the polynomial \( f = 4t^3 - 3t - \cos \theta \) in \( \mathbb{Q}(\cos \theta) \) is reducible if and only if \( f \) has a root in \( \mathbb{Q}(\cos \theta) \).

**Examples 51.13.**
1. Let \( \theta = 60^\circ \), so \( \cos \theta = 1/2 \), \( \mathbb{Q} = \mathbb{Q}(\cos \theta) \). The only possible roots of \( f = 4t^3 - 3t - 1/2 \) in \( \mathbb{Q} \) are \( \pm 1, \pm 1/2, \pm 1/4, \pm 1/8 \), none of which are, so \( f \) is irreducible in \( \mathbb{Q}[t] \), hence an angle of \( 60^\circ \) cannot be trisected with straight-edge and compass.
2. Let \( \theta = 90^\circ \), so \( \cos \theta = 0 \) and \( \mathbb{Q} = \mathbb{Q}(\cos \theta) \). As the polynomial \( f = 4t^3 - 3t \) is reducible in \( \mathbb{Q}[t] \), the angle \( 90^\circ \) can be trisected with straight-edge and compass. (Of course, we already knew this).
3. Let \( \theta = 180^\circ \), so \( \cos \theta = -1 \) and \( \mathbb{Q} = \mathbb{Q}(\cos \theta) \). As the polynomial \( f = 4t^3 - 3t + 1 \) has \(-1\) as a root, it is reducible in \( \mathbb{Q}[t] \), so the angle \( 180^\circ \) can be trisected with straight-edge and compass. (Of course, this also implies that an angle of \( 90^\circ \) can be trisected.) Note that this means that an equilateral triangle can be constructed from two given points.

**Doubling the cube:** We are given \( z_1 = 0, z_2 = 1 \) and we wish to construct a line segment of length \( a \) such that \( a^3 = 2 \) from \( z_1 = 0, z_2 = 1 \). As the polynomial \( m_\mathbb{Q}(\sqrt[3]{2}) = t^3 - 2 \) in \( \mathbb{Q}[t] \) is irreducible (by Eisenstein’s Criterion), \( \deg m_\mathbb{Q}(\sqrt[3]{2}) = 3 \), and hence this construction cannot be done with straight-edge and compass from the given two points.

**Squaring the circle:** Given a unit circle, \( z_1 = 0 \) and \( z_2 = 1 \), can we construct a segment of length \( \sqrt{\pi} \) from these two points? We can construct such a segment of this length from \( z_1 = 0 \) and \( z_2 = 1 \) if and only if we can construct a segment of length \( \pi \) from these two points and if this can be done, then \( \deg m_\mathbb{Q}(\pi) = 3 \), and hence this construction cannot be done with straight-edge and compass from the given two points.

**Construction of regular \( n \)-gons:** We must determine when it is possible to construct a ray of angle \( 2\pi/n \) radians from \( z_1 = 0, z_2 = 1, \ldots, z_{n-1} = 1 \),
equivalently, when it is possible to construct \( z = \cos \frac{2\pi}{n} + \sqrt{-1} \sin \frac{2\pi}{n} \) from \( z_1, z_2 \). To solve this construction problem, we must use the Square Root Tower Theorem (to be proven in 56.11 below). We begin with two lemmas.

**Lemma 51.14.** Any \( 2^n \)-gon, \( n > 1 \), is constructible from two points.

**Proof.** As we can erect perpendiculars, i.e., angles \( 2\pi/4 \) radians, we can construct a square. As we can bisect any angle, we can construct a ray of \( \frac{1}{2^m} \) \((\pi/2)\) for any \( m \) by induction. \( \square \)

**Lemma 51.15.** Let \( m \) and \( n \) be relatively prime integers each at least three. Then

1. A regular \((mn)\)-gon is constructible from two points if and only if both a regular \( m \)-gon and a regular \( n \)-gon can be constructible from two points.

2. A regular \( n \)-gon can be constructed from two points if and only if a regular \( 2n \)-gon can be constructed from two points.

**Proof.** (2): Angles can be bisected or doubled using straight-edge and compass.

(1): Suppose that a regular \((mn)\)-gon can be constructed, i.e., a ray of angle \( \theta = 2\pi/mn \) radians. Then rays of angle \( m\theta \) and \( n\theta \) can be constructed. Conversely, suppose that \( \theta = 2\pi/m \) radians and \( \varphi = 2\pi/n \) radians are constructible. As \( m \) and \( n \) are relatively prime, we have an equation \( 1 = am + bn \) for some integers \( a \) and \( b \). So

\[
\frac{1}{mn} = \frac{a}{n} + \frac{b}{m}
\]

hence

\[
\frac{2\pi}{mn} = \frac{2\pi}{n} + \frac{2\pi}{m}
\]

is constructible. \( \square \)

We call a number \( f_n = 2^{2^n} + 1 \) a Fermat number. If it is a prime, we call it a Fermat prime.

**Remarks 51.16.**

1. \( f_0 = 3, \ f_1 = 5, \ f_2 = 17, \ f_3 = 257, \) and \( f_4 = 65537 \) are Fermat primes.

2. (Euler) \( f_5 = 4294967297 = 641 \cdot 6700417 \).

3. \( f_6 = 274177 \cdot 67280421310722 \) is composite.

4. \( f_n \) for \( n = 5–32, \ 36–39, \ 42–43, \ 48, \ 52, \ 55, \ 58, \ 61–64, \ 66, \ 71–73, \ 75, \ 77, \ 81, \ 88, \ 90–91, \ 93–94, \ 99 \) are known to be composite.

At the time of this writing \( f_{2543548} \) is the largest composite Fermat number known and \( f_4 \) is the largest Fermat prime known.
The natural question is whether there exist infinitely many Fermat primes, or even one $f_n$, with $n > 4$.

Fermat primes are applicable to solving the construction of regular $n$-gons, because the following theorem solves this Euclidean problem.

Theorem 51.17. Let $n$ be an integer at least three. Then a regular $n$-gon can be constructed from two points if and only if $n = 2^k p_1 \cdots p_r$ with $k$ non-negative, and $p_1, \ldots, p_r$ distinct Fermat primes for some $r$ or $n = 2^k$ for some integer $k$ at least two.

Proof. By the lemmas, it suffices to determine when a regular $p^r$-gon can be constructed with $p$ an odd prime, $r \geq 1$, i.e., to determine when $ζ = \cos \frac{2π}{p^r} + \sqrt{-1} \sin \frac{2π}{p^r}$ is constructible. By Example 49.15(6), we know that $[Q(ζ) : Q] = p^r - 1$. By the Square Root Tower Theorem, $ζ$ is constructible if and only if $p^r - 1 = 2^e$ for some $e \geq 0$. For this to be true, we must have $r = 1$ and $p = 2^e + 1$ a prime. If $e = ab$, for positive integers $a$ and $b$ with $a$ odd, then

$$p = 2^{ab} + 1 = (2^b + 1)(2^{b(a-1)} - 2^{b(a-2)} + \cdots - 2^b + 1)$$

is a composite number unless $u = 1$, i.e., $p = 2^{2^a} + 1$ is a Fermat prime. This proves the theorem. □

Examples 51.18. 1. A regular 9-gon cannot be constructed either by the $n$-gon theorem or by our results on trisection of a angle. Indeed, if the 9-gon could be constructed, a ray of 40° could be constructed, hence an angle of 120° could be trisected, so an angle of 60° could be trisected which we know is false.

2. We show that a regular pentagon can be constructed, by producing the necessary square root tower. We must construct $z = e^{2π\sqrt{-1}/5}$, a root of $f = t^4 + t^3 + t^2 + t + 1$ in $Q(t)$. We have $z^{-1} = z^4 = e^{8π\sqrt{-1}/5} = e^{-2π\sqrt{-1}/5} = \overline{z}$ and $z^4 + z^3 + z^2 + z + 1 = 0$. So $z, \overline{z}, z^2, \overline{z^2}$ are the roots of $f$ in $C$. Let

$$x_1 = z + \overline{z} \text{ and } x_2 = z^2 + \overline{z^2},$$

then

$$x_1 + x_2 = -1$$

$$x_1x_2 = (z + \overline{z})(z^2 + \overline{z^2}) = -1,$$

so $(t - x_1)(t - x_2) = t^2 + t + 1$. This is not quite what we want, so we complete the square. Let $y_1 = 2x_1 + 1$ and $y_2 = 2x_2 + 1$, then

$$y_1 + y_2 = 0 \text{ and } y_1y_2 = 5.$$
Hence $t^2 - 5 = (t - y_1)(t - y_2)$ with $t^2 - 5$ irreducible in $\mathbb{Q}[t]$. We may assume that $y_1 = \sqrt{5}$. We have $x_1, x_2 \in \mathbb{Q}(y_1, y_2) = \mathbb{Q}(y_1) = \mathbb{Q}(\sqrt{5})$ and

$$
\begin{align*}
x_1^2 &= (z + \bar{z})^2 = z^2 + 2z\bar{z} + \bar{z}^2 = x_2 + 2 \\
(z - \bar{z})^2 &= z^2 - 2z\bar{z} + \bar{z}^2 = x_2 - 2.
\end{align*}
$$

Let $K = \mathbb{Q}(y_1, z - \bar{z})$. As $(z - \bar{z})^2 \in \mathbb{Q}(y_1) = \mathbb{Q}(\sqrt{5})$, the field extension $K/\mathbb{Q}$ is a square root tower. As $z = (z - \bar{z} + x_1)$ lies in $K$, a regular pentagon can be constructed.

To actually do the construction in the complex plane, draw a circle at the origin of radius two. Let $O$ be the origin, $V = \sqrt{-1}$, and $P_0 = 2$ in the plane. Let $N$ be the intersection point of the $X$-axis and the bisector of the angle $\theta = OV P_0$ and $P$ be the intersection point of the circle (in upper half plane) and the vertical constructed at $N$. Then the point $N$ is $2z = 2 \cos(2\pi/5) + 2\sqrt{-1}\sin(2\pi/5)$. (Cf. Figure 51.19.)

\begin{figure}[h]
\centering
\includegraphics[width=0.5\textwidth]{figure51.19.png}
\caption{}
\end{figure}

In the figure $\tan \theta = 1/2$, so the length of the line segment $VP_0$ is $\sqrt{5}$.

Gauss proved that one could construct the regular 17-gon in 1796 at the age of 19. This was the first important euclidean construction since ancient times. He showed

$$
16 \cos \frac{2\pi}{17} = -1 + \sqrt{17} + \sqrt{34 - 2\sqrt{17} + 2\sqrt{17 + 3\sqrt{17}}} - \sqrt{34 - 2\sqrt{17} - 2\sqrt{34 + 2\sqrt{17}}}.
$$

An actual construction was shown a few years later by Johannes Erchinger. Gauss also proved that if a regular $n$-gon was constructible then $n$ was of the form in the theorem in 1801. He also stated that this was a necessary condition, but his proof (if he had one) has never been
found. Indeed it is doubtful that Gauss actually had a valid proof, since he made this statement prior to the birth of Galois and Galois theory is needed to prove necessity. The first published proof of necessity was given by Pierre Wantzel, who also proved that you could not trisect and angle or double a cube in 1837. Few have heard of Wantzel, so the truism that solving well-known problems will bring you lasting fame is not so well-founded. A regular 257-gon was constructed by Friedrich Richelot (and Schwendenwein) in 1832. The regular 65537-gon was constructed by Johann Hermes in 1895. Apparently Hermes’ actual construction rests in a trunk in the attic at the Mathematics Institute at Göttingen (presumably studied by no one). The squaring of the circle was the last of the problems solved when Lindemann proved the transcendence of $\pi$ in 1882.

**Exercises 51.20.**

1. Prove that the equations in the last paragraph in the proof of Theorem 51.4 have a solution in $C' \cap R$.

2. Show the field of complex constructible numbers over $\mathbb{Q}$ is closed under quadratic extensions, but has extensions of every odd prime degree. Does it have extensions of any odd degree? of nonzero even degree? of any degree but two?

3. Show the construction of the regular 5-gon in Figure 51.19 works.

4. Show an angle of $2\pi/n$ radians, $n$ an integer, can be trisected by straight-edge and compass if and only if $n$ is not divisible by 3.

### 52. Separable Elements

A real polynomial that has multiple roots can cause problems as that means that such a root is also a root of its derivative, hence its graph is more complicated. If the polynomial is irreducible, this cannot happen. However, if the irreducible polynomial has coefficients in a field of positive characteristic, this can occur. In this section, we study this phenomena, formally defining some concepts that we have already used and leaving many of the details as exercises (although some will be proven later). We recall some definitions.

**Definition 52.1.** Let $K/F$ be an extension of fields, $f$ a non-constant polynomial in $F[t]$ with a root $\alpha$ in $K$. We say that $\alpha$ is a *multiple root* of $g$ of *multiplicity* $r$ in $K$ if $f = (t - \alpha)^r g$ for some $g$ in $K[t]$ with $\alpha$ not a root of $g$. If $\alpha$ is a root but not a multiple root, we call it a *simple* root of $f$. If $f = \sum_{i=0}^{n} a_i t^i$ in $F[t]$, the *derivative* of $f$ is defined to be the polynomial $f' = \sum_{i=1}^{n} i a_i t^{i-1}$ in $F[t]$. 

We have, that we (again) leave as an exercise:

**Properties 52.2.** Let $f$ and $g$ be polynomials in $F[t]$. Then

1. $(f + g)' = f' + g'$.
2. $(af)' = af'$ for all $a \in F$.
3. $(fg)' = f'g + fg'$.

As mentioned above, multiple roots are also roots of the derivative. We formally write this as a lemma. This lemma and its proof also says that the notion of multiple root and its multiplicity is independent of the field extension of the base field.

**Lemma 52.3.** Let $K/F$ be an extension of fields, $f$ a non-constant polynomial with a root $\alpha$ in $K$. Then $\alpha$ is a simple root of $f$ if and only if $\alpha$ is not a root of $f'$.

**Proof.** We can write $f = (t - \alpha)^rg$ for some $g$ in $K[t]$ with $g(\alpha) \neq 0$ for some integer $r$. Therefore, $f' = r(t - \alpha)^{r-1}g + (t - \alpha)^rg'$ in $K[t]$. As $g(\alpha) \neq 0$ in the domain $K$, we have $t - \alpha | f'$ in $K[t]$ if and only if $t - \alpha | r(t - \alpha)^{r-1}g$, i.e., $r - 1 \geq 1$ or $r = 0$ in $K$. If $r = 0$ in $K$, then $\text{char } F | r$, so $r \geq \text{char } F \geq 2$. The result follows.

**Corollary 52.4.** Let $f$ be a non-constant polynomial in $F[t]$. Then $f$ has a multiple root in some extension field $K$ of $F$ if and only if the ideal $(f, f')$ in $F[t]$ is not the unit ideal.

**Proof.** ($\Rightarrow$): Let $\alpha$ be a multiple root of $f$ in an extension field $K$ of $F$. Then $m_F(\alpha) | f$ and $m_F(\alpha) | f'$ in $F[t]$, so $(f, f') \subset (m_F(\alpha)) < F[t]$.

($\Leftarrow$): $F[t]$ is a PID, so $(f, f') = (g) < F[t]$ for some non-constant polynomial $g$ in $F[t]$. Let $h$ be any non-constant polynomial in $F[t]$ satisfying $h | g$. Let $K/F$ be a field extension such that $h$ has a root $\alpha$ in $K$. Then $f(\alpha) = f'(\alpha) = h(\alpha) = 0$, so $f$ has a multiple root in $K$.

This also answers the question of multiple roots for irreducible polynomials.

**Corollary 52.5.** Let $f$ be an irreducible polynomial in $F[t]$.

1. If char $F = 0$, then $f$ has no multiple roots in any field extension of $F$. In particular, $f$ can have only simple roots, if any, in a field extension of $F$.

2. If char $F = p > 0$ and $F$ has a multiple root in some field extension $K$ of $F$, then there exists a polynomial $g$ in $F[t]$ satisfying $f = g(t^p)$, i.e., $f = \sum a_it^{pi}$ in $F[t]$. 
Proof. We begin the proof, leaving the rest of the proof as exercises. Let $K/F$ be a field extension such that $f$ has a root $\alpha$ in $K$, so $f = am_F(\alpha)$ for some nonzero $a$ in $F$. In particular, if $\alpha$ is a multiple root, we must have $f \mid f'$ in $F[t]$. So $f' = 0$ or $\deg f' \geq \deg f$. Consequently, $f' = 0$. The result now follows by Exercises 52.10(2a) and (2b). □

Definition 52.6. An irreducible polynomial $f$ in $F[t]$ is called separable over $F$ if it has no multiple roots in any field extension of $F$. A non-constant polynomial is called separable over $F$ if all of its irreducible factors are separable. If $f$ is not separable over $F$, it is called inseparable. Let $K/F$ be a field extension and $\alpha$ an element in $K$. We say that $\alpha$ is separable over $F$ if it is the root of a separable polynomial defined over $F$ and the extension $K/F$ is separable if every element of $K$ is separable over $F$.

Remark 52.7. In reading other books, it is much more common to say that a non-constant polynomial $f$ in $F[t]$ is separable if it has no multiple roots in any extension field of $F$. It is, however, convenient to use the above definition for a separable polynomial in the sequel.

Remarks 52.8. Let $K/F$ be an extension of fields and $f$, $g$ non-constant polynomials in $F[t]$.

1. Every element of $F$ is separable over $F$.
2. An element $\alpha$ in $K$ is separable over $F$ if and only if $\alpha$ is algebraic over $F$ and $m_F(\alpha)$ is separable over $F$. In particular, if $K/F$ is separable, then it is algebraic.
3. If char $F = 0$ and $K/F$ is algebraic, then $K/F$ is separable. A field $F$ in which every element algebraic over $F$ is separable is called a perfect field. Therefore, any field of characteristic zero is perfect. It turns out that every finite field is also perfect, but there exist fields that are not perfect.
4. The polynomial $f$ is separable over $F$ if and only if every divisor of $f$ is separable over $F$.
5. If both $f$ and $g$ are separable over $F$ then so is $fg$.
6. If $L/K/F$ are field extensions with $L/F$ separable, then $L/K$ and $K/F$ are separable.

Theorem 52.9. Let $K/E/F$ be extensions of field. Then $K/F$ is separable if and only if $K/E$ and $E/F$ are both separable.

We shall not prove this theorem here. We know that if $K/F$ is separable, then $K/E$ and $E/F$ are separable. The converse is not so
easy. To reduce the problem, we can try a trick that we did before. If $\alpha$ in $K$ is separable over $E$, then $m_E(\alpha)$ in $E[t]$ is separable. If $m_E(\alpha) = t^n + a_{n-1}t^{n-1} + \cdots + a_0$ in $E[t]$, then $m_E(\alpha)$ is separable over $E_0 = F(a_0, \ldots, a_{n-1})$, with $a_i$, $i = 0, \ldots, n - 1$ separable over $F$.

So we are reduced to showing if $\beta_1, \ldots, \beta_n$ are separable over $F$, then $F(\beta_1, \ldots, \beta_n)/F$ is separable and if, in addition, $\alpha$ is separable over $F(\beta_1, \ldots, \beta_n)$, then $\alpha$ is separable over $F$. This is still not easy. One way of doing this is given in the exercises. Of course, we may assume that char $F$ is positive.

Exercises 52.10.

1. Suppose that $f$ be a non-constant polynomial in $F[t]$, $L/K_i/F$ field extensions with $i = 1, 2$, and $\alpha$ an element of $K_1 \cap K_2$. Show that $\alpha$ is a root of $f$ of multiplicity $r$ in $K_1$ if and only if $\alpha$ is a root of $f$ of multiplicity $r$ in $K_2$.

2. Let $F$ be a field and $f$ a polynomial in $F[t]$. Show the following:
   (a) If char $F = 0$, then $f' = 0$ if and only if $f$ is a constant polynomial.
   (b) If char $F = p$ is positive, then $f(t)^p = f^p(t^p)$ (The Children’s Binomial Theorem) and if $f' = 0$, then there exists a polynomial $g$ in $F[t]$ satisfying $f(t) = g(t^p)$.

3. If $x$ is transcendental over a field $F$, then $t^p - x \in F(x)[t]$ is irreducible for any prime $p$.

4. Suppose that $F$ is a field of positive characteristic $p$. Show the following:
   (a) The map $F \to F$ given by $x \mapsto x^p$ is a monomorphism. Denote its image by $F_p$.
   (b) If $K/F$ is algebraic and $\alpha \in K$ is separable over $F(\alpha^p)$, then $\alpha \in F(\alpha^p)$.
   (c) Every finite field is perfect, i.e., every algebraic extension is separable.

5. Suppose that $F$ is a field of positive characteristic $p$. Show all of the following:
   (a) If $K/F$ is separable field extension, then $K = F(K^p)$, where $K^p := \{x^p \mid x \in K\}$.
   (b) Suppose that $K/F$ is finite and $K = F(K^p)$. If $\{x_1, \ldots, x_n\} \subset K$ is linearly independent over $F$ then so is $\{x_1^p, \ldots, x_n^p\}$.
   (c) If $K/F$ is finite and $K = F(K^p)$, then $K/F$ is separable.

6. Let $K/F$ be an extension of fields. Show all of the following:
   (a) If $\alpha \in K$ is separable over $F$ then $F(\alpha)/F$ is separable.
52. SEPARABLE ELEMENTS

(b) If \( \alpha_1, \ldots, \alpha_n \in K \) are separable over \( F \), then \( F(\alpha_1, \ldots, \alpha_n)/F \) is separable.

(c) Let \( F_{\text{sep}} = \{ \alpha \in K \mid \alpha \text{ separable over } F \} \). Then \( F_{\text{sep}} \) is a field.

7. Let \( L/K/F \) be field extensions. Show that \( L/F \) is separable if and only if \( L/K \) and \( K/F \) are separable using the previous two exercises.

8. Show that any algebraic extension of a perfect field is perfect.

9. Let \( F_0 \) be a field of positive characteristic \( p \). Let \( F = F_0(t_1^p, t_2^p) \) and \( L = F_0(t_1, t_2) \). Show all of the following:
   (a) If \( \theta \in L \setminus F \) then \( [F(\theta) : F] = p \).
   (b) There exist infinitely many fields \( K \) satisfying \( F < K < L \).

10. Let \( F \) be a field of positive characteristic \( p \). Show
   (a) If \( F = F^p \), then \( F \) is perfect.
   (b) If \( F \) is not perfect, then \( F \neq F^p := \{ x^p \mid x \in F \} \) for any \( r \geq 1 \).

11. Let \( K/F \) be a field extension and \( \alpha \in K \). We say that \( \alpha \) is \textit{purely inseparable} over \( F \) if \( m_F(\alpha) \) has only one root in a splitting field of \( \alpha \) and \( K/F \) is \textit{purely inseparable} if every element of \( K \) is purely inseparable over \( F \). Show that \( \alpha \) in \( K \) is pure inseparable and separable over \( K \) if and only if \( \alpha \) lies in \( F \).

12. Let \( K/F \) be a field extension with \( \text{char}(F) = p > 0 \). Show that if \( \alpha \) in \( K \) is algebraic then there exists a nonnegative integer \( n \) satisfying \( \alpha^{p^n} \) is separable over \( F \).

13. Let \( K/F \) be an algebraic extension with \( \text{char}(F) = p > 0 \). Show the following are equivalent:
   (a) \( K/F \) is purely inseparable.
   (b) If \( \alpha \) in \( K \), then there exists a nonnegative integer \( n \) such that \( m_F(\alpha) = t^{p^n} - a \in F[t] \).
   (c) If \( \alpha \) lies in \( K \), then there exists a nonnegative integer \( n \) such that \( \alpha^{p^n} \) lies in \( F \).
   (d) \( K_{\text{sep}} = F \). (Cf. Exercise 6 above.)
   (e) \( K \) is generated by purely inseparable elements over \( F \).

14. Let \( K/F \) be an algebraic field extension. Show that \( K/F_{\text{sep}} \) is purely inseparable and \( K_{\text{psep}} := \{ \alpha \in K \mid \alpha \text{ is purely inseparable over } F \} \) is a field.
CHAPTER XII

Galois Theory

In the previous chapter, we studied field theory, i.e., fields and field extensions. In the chapter, we interrelate field theory with group theory. The goal will be to intertwine the intermediate fields of finite field extension $K/F$ when $K$ is the splitting field of a separable polynomial in $F[t]$ over $F$ and subgroups of the Galois group $\text{Aut}_F(K)$. This will allow us to prove some fundamental results. For example, we shall give an algebraic proof of the Fundamental Theorem of Algebra using only the Intermediate Value Theorem from analysis, prove that there is no formula involving only addition, multiplication, and the extraction of $n$th roots for various $n$ for the fifth degree polynomial over the rational numbers using the results that were established about solvable groups, and finish the proof for the existence of the regular $n$-gon for the allowable $n$. We shall show how the study of roots of unity leads to Quadratic Reciprocity and a special case of Dirichlet’s Theorem of Primes in an Arithmetic Progression.

53. Characters

This section is the key to the interaction of field theory and group theory. We follow the approach of E. Artin, who proved the main results using systems of linear equations. The idea is to look at a set $S$ of field homomorphisms from a field $K$ and to determine the subfield of $K$ on which each homomorphism in $S$ acts in the same way. We begin with the following definition:

Definition 53.1. Let $F$ be a field and $G$ be a (multiplicative) group. Then a group homomorphism $\sigma : G \to F^\times$ is called an (linear) character. [Note that $F^\times = \text{GL}_1(F)$.] We say that distinct characters $\sigma_1, \ldots, \sigma_n : G \to F^\times$ are dependent (or the set $\{\sigma_1, \ldots, \sigma_n\}$ is dependent) if there exist $a_1, \ldots, a_n$ in $F$, not all zero, satisfying $\sum a_i\sigma_i(g) = 0$ for all $g$ in $G$, i.e., the function $\sum a_i\sigma_i : G \to F$ is the zero map; and independent otherwise.

Remarks 53.2. 1. A field homomorphism $\sigma : E \to F$ induces an character by restriction, as a ring homomorphism takes units to
units. As field homomorphisms are monic, such a character is always injective.

2. The set

\[ \{ \sigma \mid \sigma : G \to F^\times \text{ is a character.} \} \]

is not a group under +, but is a group under multiplication of functions, i.e., \( \sigma \tau (g) := \sigma (g) \tau (g) \) for all \( g \) in \( G \), with the identity the map \( 1 : G \to F^\times \) given by \( g \mapsto 1_F \).

The key result is:

**Lemma 53.3.** (Dedekind’s Lemma) Let \( \sigma_1, \ldots, \sigma_n : G \to F^\times \) be distinct characters. Then \( \sigma_1, \ldots, \sigma_n : G \to F^\times \) are independent.

**Proof.** We prove this by induction on \( n \).

\( n = 1 \): If \( a \sigma_1 (g) = 0 \) for all \( g \) in \( G \), then \( a = 0 \) as \( \sigma_1 (g) \neq 0 \) for all \( g \) in \( G \).

\( n > 1 \): Suppose that we have an equation

\[ (*) \quad \sum_{i=1}^{n} a_i \sigma_i (g) = 0 \text{ for all } g \text{ in } G \text{ with } a_i \text{ in } F \text{ not all zero.} \]

By induction any proper subset of \( \{ \sigma_1, \ldots, \sigma_n \} \) is independent, so \( a_i \) is nonzero for every \( i \). Multiplying the equation in \( (*) \) by \( a_n^{-1} \), we may assume that that \( a_n = 1 \). So \( (*) \) now reads

\[ (**) \quad \sum_{i=1}^{n-1} a_i \sigma_i (g) + \sigma_n (g) = 0 \text{ for all } g \text{ in } G. \]

By assumption, \( \sigma_1 \neq \sigma_n \), so there exists an element \( x \) in \( G \) satisfying \( \sigma_1 (x) \neq \sigma_n (x) \). As \( G \) is a group, \( xG = G \); and as the \( \sigma_i \) are group homomorphisms, we have

\[ 0 = \sum_{i=1}^{n-1} a_i \sigma_i (xg) + \sigma_n (xg) = \sum_{i=1}^{n-1} a_i \sigma_i (x) \sigma_i (g) + \sigma_n (x) \sigma_n (g) \]

for all \( g \) in \( G \). Since \( \sigma_n (x) \neq 0 \), we see that

\[ (\dagger) \quad 0 = \sum_{i=1}^{n-1} \sigma_n (x)^{-1} a_i \sigma_i (x) \sigma_i (g) + \sigma_n (g) \text{ for all } g \text{ in } G. \]

Subtracting the equation in \( (\dagger) \) from the equation in \( (**) \) yields

\[ 0 = \sum_{i=1}^{n-1} [a_i - \sigma_n (x)^{-1} a_i \sigma_i (x)] \sigma_i (g) \text{ for all } g \text{ in } G. \]
By induction, we conclude that
\[ a_i = \sigma_n(x)^{-1}a_i\sigma_i(x) \] for all \( i \)
in the domain \( F \) with \( a_i \neq 0 \), hence \( \sigma_n(x) = \sigma_i(x) \) for all \( i \), contradicting \( \sigma_1(x) \neq \sigma_n(x) \).

The proof above also shows, with the obvious definitions, that any set of characters \( S = \{ \sigma_i : G \to F^\times \mid \sigma_i \text{ a character for all } i \in I \} \) is independent.

**Corollary 53.4.** Let \( \sigma_1, \ldots, \sigma_n : F \to K \) be distinct field homomorphisms. Then \( \sigma_1, \ldots, \sigma_n : F^\times \to K^\times \) are independent characters.

**Definition 53.5.** Let \( S \) be a non-empty subset of the full set of all field homomorphisms \( F \to K \), \( \{ \sigma : F \to K \mid \sigma \text{ a field homomorphism} \} \). We call an element \( a \) in \( F \) a \textit{fixed point under} \( S \), if \( \sigma(a) = \tau(a) \) for all \( \sigma \) and \( \tau \) in \( S \).

**Example 53.6.** Let \( S \) be a non-empty set of field homomorphisms \( F \to K \). Assume that both \( F \) and \( K \) contain the same prime subfield \( \Delta \) (so \( \Delta \cong Q \) if the characteristic of \( F \) is zero and \( \Delta \cong \mathbb{Z}/p\mathbb{Z} \), if the characteristic of \( F \) is the prime \( p \)). As every homomorphism takes \( 1_F \to 1_K \) and \( 1_{\Delta} = 1_F = 1_K \), we have \( \sigma(x) = x \) for all \( x \in \Delta \).

Therefore, the term ‘fixed point’ is actually appropriate for elements in \( \Delta \).

Note that if \( S \) is a set of field homomorphisms \( \sigma : K \to K \) containing \( 1_K \), then the fixes points of this set are all really fixed points, i.e., \( \sigma(x) = x \) for all \( \sigma \in S \).

As should be expected, the set of fixed points has additional structure. We immediately see that we have:

**Lemma 53.7.** Suppose that \( S \) is a non-empty set of field homomorphisms \( F \to K \). Let \( E = \{ x \in F \mid x \text{ is a fixed point of } S \} \). Then \( E \) is a subfield of \( F \).

If \( S \) is a non-empty set of field homomorphisms \( F \to K \), we denote the field \( \{ x \in F \mid x \text{ is a fixed point of } S \} \) by \( F^S \) and call it the \textit{fixed field} of \( S \). We are interested in the fixed field of a finite set of field homomorphisms. Following Artin, we use linear algebra to obtain non-trivial solutions to an appropriate set of linear equations. The first main result about such is the following:

**Lemma 53.8.** (Artin’s Lemma) Suppose that \( S \) is a finite non-empty set of field homomorphisms \( F \to K \). Then \( [F : F^S] \geq |S| \).
Proof. Let $S = \{\sigma_1, \ldots, \sigma_n\}$ with $n \geq 1$. Suppose that $r = [F : F^S] < |S| = n$ and $\{\omega_1, \ldots, \omega_r\}$ is an $F^S$-basis for $F$. Consider the following system of equations over $K$:

\[
\begin{align*}
\sigma_1(\omega_1)x_1 + \cdots + \sigma_n(\omega_1)x_n &= 0 \\
\vdots \\
\sigma_1(\omega_r)x_1 + \cdots + \sigma_n(\omega_r)x_n &= 0.
\end{align*}
\]

As $r < n$, there exists a non-trivial solution to the system of equations (*) over $K$, say $x_1, \ldots, x_n$ with each $x_i$ in $K$ and not all zero is such a solution. Let $a$ be an element of $F$. We can write $a = \sum a_i \omega_i$ with each $a_i \in F^S$. For each $i$, multiply the $i$th equation in (*) by $\sigma_i(a_i)$ to get a new equivalent system of equations

\[
\begin{align*}
\sigma_1(a_1)\sigma_1(\omega_1)x_1 + \cdots + \sigma_n(a_1)\sigma_n(\omega_1)x_n &= 0 \\
\vdots \\
\sigma_1(a_r)\sigma_1(\omega_r)x_1 + \cdots + \sigma_n(a_r)\sigma_n(\omega_r)x_n &= 0.
\end{align*}
\]

Since $\sigma_i(a_i) = \sigma_j(a_i)$ for all $i$ and $j$ and the $\sigma_i$ are homomorphisms, we have

\[
\begin{align*}
\sigma_1(a_1\omega_1)x_1 + \cdots + \sigma_n(a_1\omega_1)x_n &= 0 \\
\vdots \\
\sigma_1(a_r\omega_r)x_1 + \cdots + \sigma_n(a_r\omega_r)x_n &= 0.
\end{align*}
\]

Adding all the equations in (†) yields

\[
0 = \sum_{i=1}^{r} \sum_{j=1}^{n} \sigma_j(a_i \omega_i) x_j = \sum_{j=1}^{n} \sum_{i=1}^{r} \sigma_j(a_i \omega_i) x_j = \sum_{j=1}^{n} \sigma_j(a) x_j
\]

for all $a$ in $F$ with $x_j$ not all zero. This means that the characters $\sigma_1, \ldots, \sigma_n$ are dependent, contradicting Dedekind’s Lemma. $\square$

A special case of Artin’s Lemma, is the one in which we shall be interested.

Corollary 53.9. Let $S$ be a non-empty finite set of field automorphisms of a field $K$. Then $[K : K^S] \geq |S|$.

Of course, we should be very interested in the case when we have equality in Artin’s Lemma. In general this is not true. We can even have a finite set of field automorphisms $S$ of a field $K$ with $[K : K^S]$ infinite [can you give an example?]. What we need is to look at those
S with additional structure. For example, let $K/F$ be an extension of fields. Recall the Galois group of $K/F$ is defined to be the set

$$G(K/F) = \text{Aut}_F(K)$$

$$ := \{ \sigma : K \to K \mid \text{an } F\text{-automorphism satisfying } \sigma|_F = 1_F \}.$$ 

It is immediate that the following is true.

**Lemma 53.10.** Let $K/F$ be an extension of fields. Then $G(K/F)$ is a subgroup of the (field) automorphism group of $K$ and $F$ is a subfield of $K^{G(K/F)}$.

As a consequence, we have:

**Corollary 53.11.** Let $K/F$ be a finite extension of fields. Then the Galois group $G(K/F)$ is a finite group and satisfies $[K : F] \geq |G(K/F)|$.

**Proof.** As $F$ is a subfield of $K^S$ for any non-empty subset $S$ of $G(K/F)$, we have

$$[K : F] \geq [K : K^S] \geq |S|$$

for any finite $S$.

We arrive at the group theoretic condition in which we are interested.

**Definition 53.12.** Let $K/F$ be a finite extension of fields. We say that the extension is Galois if $F$ is the fixed field of $G(K/F)$, i.e.,

$$F = K^{G(K/F)}.$$ 

We have defined a field extension $K/F$ to be Galois if $K$ is a finite extension with $F$ the fixed field of $G(K/F)$. It would be more appropriate to call this a finite Galois extension and call $K/F$ a Galois extension if it is algebraic with $F$ the fixed field of $G(K/F)$.

**Remarks 53.13.** Let $K/F$ be a field extension.

1. If $F = K^{G(K/F)}$, then for any element $x \in K \setminus F$, there exists an automorphism in $G(K/F)$ that moves $x$, i.e., there exists a $\sigma$ in $G(K/F)$ satisfying $\sigma(x) \neq x$.

2. If $K/E/F$ is an intermediate field, then $G(K/E) \subset G(K/F)$.

3. Let $E = K^{G(K/F)}$, then by definition of $K/E/F$, we have $G(K/E) \supset G(K/F)$, so by (2), we have

$$G(K/F) = G(K/K^{G(K/F)}).$$ 

4. If $K/F$ is finite, then a field homomorphism $\sigma : K \to K$ fixing $F$ necessarily is an element of $G(K/F)$ as $\sigma$ is monic and $F$-linear.
We give examples of Galois groups.

**Examples 53.14.** 1. Let $K = F$ be a field. Then $G(K/F) = \{1\}$ and $K/F$ is Galois with $[K : F] = 1 = |G(K/F)|$. Moreover, $K$ is the splitting field of $t - a$ in $F[t]$ over $F$ for any $a \in F$.

2. Let $F = \mathbb{R}$ and $K = \mathbb{C}$. Then $G(K/F) = \{1, -\}$ (with $1 = 1_K$ and $-$ complex conjugation) and $K/F$ is Galois with $[K : F] = 2 = |G(K/F)|$. Moreover, $K$ is the splitting field of $t^2 + 1$ over $F$.

3. Suppose that $F = \mathbb{Q}$ and $K = \mathbb{Q}(\sqrt{2})$. Let $\sigma : K \to K$ be defined by $a + b\sqrt{2} \mapsto a - b\sqrt{2}$ for all $a$ and $b$ in $\mathbb{Q}$. Then $G(K/F) = \{1, \sigma\}$ (with $1 = 1_K$) and $K/F$ is Galois with $[K : F] = 2 = |G(K/F)|$. Moreover, $K$ is the splitting field of $t^2 - 2$ over $F$.

4. Suppose that $F = \mathbb{Q}$ and $K = \mathbb{Q}(\sqrt[3]{2})$. Then $G(K/F) = \{1\}$, so $K/F$ is not Galois. It satisfies $[K : F] = 3 > |G(K/F)|$. Moreover, $K$ is not a splitting field of any polynomial over $F$.

We indicate why these facts are true. Any field automorphism of $K$ must take $\sqrt[3]{2}$ to another root of $t^3 - 2$, but $\mathbb{R}$ contains only one root of $t^3 - 2$. If $K$ was the splitting field of some polynomial $g$ in $F[t]$ over $F$, then it would be the splitting field of an irreducible polynomial of degree three. (Why?) Such a $g$ would have three distinct roots, say $\alpha_1, \alpha_2, \alpha_3$ in $K$ and we would have $K = \mathbb{Q}(\alpha_i)$, $i = 1, 2, 3$. But then there must be an $F$-automorphism $\sigma$ satisfying $\sigma(\alpha_1) = \alpha_2$.

5. Suppose that $F = \mathbb{Q}$ and $K = \mathbb{R}$. Then $G(K/F) = \{1\}$ and $K^{G(K/F)} = K$ with $K/F$ is infinite and not even algebraic. (Cf. Exercise 49.18(11).) However, (by Exercise 55.22(6)) the Galois group $G(\mathbb{C}/\mathbb{Q})$ is uncountable!

6. Suppose that the characteristic of a field $F$ is $p > 0$ and $K = F(\alpha)$ for some $\alpha$ satisfying $\alpha^p = a$ with $a \in F$. If $\alpha \notin F$, then $G(K/F) = \{1\}$ and $K/F$ is not Galois. It satisfies $[K : F] = p > |G(K/F)|$ and $K$ is the splitting field of $t^p - a$ over $F$.

Artin’s Lemma becomes much stronger, when the non-empty finite set of field automorphisms is a group.

**Theorem 53.15.** (Artin’s Theorem) Let $K$ be a field and $G$ a finite subgroup of field automorphisms of $K$. Then

2. $G = G(K/K^G)$.
3. $K/K^G$ is Galois.

**Proof.** (1): Let $n = |G|$ and $G = \{\sigma_1, \ldots, \sigma_n\}$. By Artin’s Lemma, we know that $[K : K^G] \geq |G| = n$. Suppose that $[K : K^G] >
n. Choose \( a_1, \ldots, a_{n+1} \) in \( K^\times \) that are \( K^G \)-linearly independent. The system of \( n \) linear equations in \((n+1)\)-unknowns

\[
\begin{align*}
x_1\sigma_1(a_1) + \cdots + x_{n+1}\sigma_1(a_{n+1}) &= 0 \\
\vdots \\
x_1\sigma_n(a_1) + \cdots + x_{n+1}\sigma_n(a_{n+1}) &= 0
\end{align*}
\]

(*) has a non-trivial solution over \( K \). Among all the non-trivial solutions \( x_1, \ldots, x_{n+1} \) over \( K \), choose one with the least number of \( x_i \) nonzero. Relabeling, we may assume this solution is

\[
x_1, \ldots, x_r, 0, \ldots, 0 \text{ with } x_i \in K^\times, \ 1 \leq i \leq r
\]

with \( r \) minimal. Multiplying this solution by \( x_r^{-1} \), we may also assume that \( x_r = 1 \). As \( G \) is a group, we may assume also that \( \sigma_1 = 1_K \), so

\[
x_1a_1 + \cdots + x_{r-1}a_{r-1} + a_r = 0, \ x_i \in K^\times, \ r > 1 \ (\text{as } a_1 \neq 0).
\]

Since \( \{a_1, \ldots, a_r, a_{r+1}, \ldots, a_{n+1}\} \) is \( K^G \)-linearly independent, this means that for some integer \( i, 1 \leq i \leq r \), we have \( x_i \notin K^G \). By definition, there exists an integer \( k \) such that \( \sigma_k(x_i) \neq x_i \). (So \( k > 1 \).) The system of equations in (*) has the form

\[
0 = \sum_{i=1}^{r-1} x_i\sigma_j(a_i) + \sigma_j(a_r) = 0 \text{ for } j = 1, \ldots, n.
\]

(**)

Taking \( \sigma_k \) of each of the equations in (**) yields

\[
0 = \sum_{i=1}^{r-1} \sigma_k(x_i)\sigma_k\sigma_j(a_i) + \sigma_k\sigma_j(a_r) = 0 \text{ for } j = 1, \ldots, n.
\]

As \( G \) is a group, \( \sigma_k G = G \), hence these equations are the same as

\[
(\dagger) \quad 0 = \sum_{i=1}^{r-1} \sigma_k(x_i)\sigma_j(a_i) + \sigma_j(a_r) = 0 \text{ for } j = 1, \ldots, n.
\]

Subtracting the appropriate equations in (\dagger) from those in (*) yields

\[
0 = \sum_{i=1}^{r-1} [x_i - \sigma_k(x_i)]\sigma_j(a_i) = 0 \text{ for } j = 1, \ldots, n.
\]

Since \( \sigma_k(x_i) \neq x_i \), we have

\[
x_1 - \sigma_k(x_1), \ldots, x_{r-1} - \sigma_k(x_{r-1}), 0, \ldots, 0
\]

is a non-trivial solution over \( K \) to the system of the equations in (*). This contradicts the minimality of \( r \) and establishes (1).
(2): By (1), we have \([K : K^G] = |G|\) is finite. By Corollary 53.9, we have \([K : K^G] \geq |G(K/K^G)|\), so \(G(K/F)\) is also finite. Since \(G\) is a subgroup of \(G(K/K^G)\) by definition, we have \(|G| = |G(K/K^G)|\) and is finite, so \(G = G(K/K^G)\).

(3) is now immediate. \(\square\)

We want to interrelate Galois groups and field extensions. An application of Artin’s Theorem produces one such result.

**Corollary 53.16.** Let \(K\) be a field, \(G_1, G_2\) two finite subgroups of \(\text{Aut}(K)\), the group of field automorphisms of \(K\), and \(F_i = K^{G_i}\) for \(i = 1, 2\). Then \(F_1 = F_2\) if and only if \(G_1 = G_2\).

**Proof.** \((\Leftarrow)\) is clear.

\((\Rightarrow)\): By Artin’s Theorem, \(G_i = G(K/F_i)\), so this also follows. \(\square\)

We now see that for finite field extensions, being Galois only depends on a simple cardinality test.

**Corollary 53.17.** Suppose that \(K/F\) is a finite extension of fields. Then \(K/F\) is Galois if and only if \([K : F] = |G(K/F)|\).

**Proof.** \((\Rightarrow)\) follows from Artin’s Theorem and the special case of Artin’s Lemma, Corollary 53.9.

\((\Leftarrow)\): Let \(E = K^{G(K/F)}\). Then we know that \(G(K/E) = G(K/F)\). As \(K/E\) is Galois, by the proven sufficiency, we have \([K : E] = |G(K/E)|\). Therefore,

\([K : E] = |G(K/E)| = |G(K/F)| = [K : E][E : F]\.

It follows that \([E : F] = 1\), i.e., \(E = F\). \(\square\)

**Exercise 53.18.** This exercise computes the Galois group of a non-algebraic extension. Using Exercise 47.24(10) show that \(G(F(t)/F)\) consists of all \(F\)-automorphisms of \(F(t)\) mapping \(t\) to \((at + b)/(ct + d)\) where \(a, b, c, d \in F\) satisfies \(ad - bc \neq 0\)

54. Computations

Of course, given a finite extension \(K/F\) of fields, one would like to determine its Galois group. In general, as expected, this is very difficult. In this section, we do a few computations as well as giving some ideas on how to do such computations. To begin, we discuss the problem by reviewing (and repeating) some of the material that we have already discussed.

**Summary 54.1.** Let \(K/F\) be a finite extension of fields and \(f\) a non-constant polynomial in \(F[t]\).
1. We know if $K/E/F$ is an intermediate field that
   (a) $G(K/E) \subset G(K/F)$.
   (b) $G(K/F) = G(K/K^{G(K/F)})$.
   (c) $[K : F] \geq |G(K/F)|$ with equality if and only if $K/F$ is Galois, 
       i.e., $F = K^{G(K/F)}$.
2. Let $\alpha$ in $K$ be a root of $f$ and $\sigma$ an element in $G(K/F)$. Then $\sigma(\alpha)$ is also a root of $f$. In particular, if $f = g_1 \cdots g_r$ in $F[t]$ with $g_i$ non-constant polynomials in $F[t]$, then $G(K/F)$ takes roots of $g_i$ in $K$ to roots of $g_i$ for all $i$:
   If $f = \sum a_i t^i$, then
   \[ 0 = \sigma(f(\alpha)) = \sigma(\sum a_i \alpha^i) = \sum \sigma(a_i) \sigma(\alpha^i) = \sum a_i \sigma(\alpha)^i = f(\sigma(\alpha)). \]
3. If $K$ is a splitting field of $f$ and $\sigma$ an element in $G(K/F)$, then $\sigma$ is completely determined by what it does to the roots of $f$:
   If $\alpha_1, \ldots, \alpha_r$ are the roots of $f$, then $K = F(\alpha_1, \ldots, \alpha_n)$.
4. Let $S = \{ \alpha \mid \alpha$ is a root of $f \}$. If $S$ is not empty, then the map $G \to \Sigma(S)$ given by $\sigma \mapsto \sigma|_S$ is a group homomorphism. If $K$ is a splitting field of $f$ over $F$, then this map is monic and we can view $G \subset \Sigma(S)$. [This was Galois’ viewpoint.]
5. If $\alpha$ is a root of $f$ in $F$ then $\sigma(\alpha) = \alpha$ for all $\sigma \in G(K/F)$.
6. This is the most useful remark. If $K$ is the splitting field of $f$ over $F$ and $g$ is an irreducible polynomial in $F[t]$ with roots $\alpha$ and $\beta$ in $K$, then there exists an $F$-isomorphism
   \[ \tau : F(\alpha) \to F(\beta) \] 
   satisfying $\alpha \mapsto \beta$
   and this extends to an $F$-automorphism $\tau$ in $G(K/F)$, so $\tau(\alpha) = \beta$.
   [Make sure that you can prove this.] In particular,
   $G(K/F)$ acts transitively on the roots of $g$ in $K$.
   We shall later see that, in fact,
   if $g$ is irreducible in $F[t]$ and 
   has a root in a splitting field $K$ of $f$, then $g$ splits over $K$.
   Can you prove this? [You have all the tools to do so.]

To do some of our computations, we shall need a few lemmas.

**Lemma 54.2.** Let $G$ be a finite cyclic group of order $n$ and $\text{Aut}(G)$ the automorphism group of $G$. Then $\text{Aut}(G) \cong (\mathbb{Z}/n\mathbb{Z})^\times$. In particular, $\text{Aut}(G)$ is abelian (and cyclic if $n$ is a prime) of order $\varphi(n)$. 
Proof. Let $G = \langle a \rangle$ and $\sigma \in \text{Aut}(G)$. Then $\sigma(a) = a^i$, for some $1 \leq i \leq n$ and $\langle a \rangle = \langle a^i \rangle$. It follows that $i$ and $n$ are relatively prime. Conversely, if $i$ and $n$ are relatively prime with $1 \leq i \leq n$, then $\sigma_i : G \to G$ given by $a^m \mapsto a^{mi}$ is checked to be a group monomorphism, hence an isomorphism as $G$ is finite. The map $(\mathbb{Z}/n\mathbb{Z})^\times \to \text{Aut}(G)$ given by $i \mod n \mapsto \sigma_i$ is checked to be an isomorphism, so we also have $|\text{Aut}(G)| = |(\mathbb{Z}/n\mathbb{Z})^\times| = \varphi(n)$. □

Remark 54.3. In E, we shall show the following: If $p$ is an odd prime then $(\mathbb{Z}/p^n\mathbb{Z})^\times$ is cyclic. If $p = 2$, this is not true. Indeed if $m \geq 3$, then $(\mathbb{Z}/2^m\mathbb{Z})^\times \cong (\mathbb{Z}/2\mathbb{Z}) \times (\mathbb{Z}/2^{m-2}\mathbb{Z})$. It then follows, using the Chinese Remainder Theorem, that $(\mathbb{Z}/n\mathbb{Z})^\times$ is cyclic if and only if $n = 2, 4, p^i$, or $2p^r$ where $p$ is an odd prime.

Proposition 54.4. Let $K$ be a splitting field of $t^n - 1$ in $F[t]$ over $F$ and

$$U = \{z \in K \mid z^n = 1\} = \{z \mid z \text{ a root of } t^n - 1 \text{ in } K\}.$$ 

Then $U$ is a cyclic subgroup of $K^\times$. Suppose, in addition, that either $\text{char } F = 0$ or $\text{char } F \nmid n$. Then $|U| = n$ and the Galois group $G(K/F)$ is isomorphic to a subgroup of $(\mathbb{Z}/n\mathbb{Z})^\times$. In particular, $G(K/F)$ is abelian and $|G(K/F)| = \varphi(n)$.

Proof. As $(z_i z_j)^n = 1 = (z_i^{-1})^n$, for all $z_i, z_j \in U$, we know that $U$ is a group. Since it is a finite subgroup of $K^\times$, it is cyclic by Theorem 33.15. As $K = F(U)$, every $\sigma$ in $G(K/F)$ is determined by $\sigma(z)$, with $z \in U$, so

$$\varphi : G(K/F) \to \Sigma(U) \text{ given by } \sigma \mapsto \sigma|_U$$

is a group monomorphism. As $U$ is cyclic, say $U = \langle z \rangle$, and every $\sigma$ in $G(K/F)$ satisfies $\sigma(z^i) = \sigma(z)^i$, it follows that every such $\sigma$ is completely determined by $\sigma(z)$ and $\sigma|_U$ lies in the automorphism group $\text{Aut}(U)$ of $U$. If $\text{char } F = 0$ or $\text{char } F \nmid n$, we have $(t^n - 1)' = nt^{n-1} \neq 0$ has no root in common with $t^n - 1$ in $F[t]$, i.e., $t^n - 1$ has only simple roots in every field extension of $F$. In particular, $|U| = n$. □

In the proposition, the element $z$ satisfying $U = \langle z \rangle$ is called a primitive $n$th root of unity. If, in the proposition, $\text{char } F = 0$ or $\text{char } F \nmid n$, then there exists $\varphi(n)$ primitive roots of unity.

Remark 54.5. A (deep) theorem of Kronecker-Weber says: Suppose a finite extension of fields $K/Q$ is Galois with abelian Galois group. We call such an extension is an abelian extension of $Q$. Then there exists a primitive $n$th root of unity $\zeta$, e.g., $e^{2\pi \sqrt{-1}/n}$, for some $n$, such that $K$ is a subfield of $Q(\zeta)$. One says that all abelian extensions of $Q$ are
cyclotomic. We shall show later this is true in the special case that \( K/F \) is quadratic.

**Remarks 54.6.** The following are left as exercises:

1. If \( z \in F \), then \( |G(K/F)| = 1 = [K : F] \).
2. Suppose that \( |F| = q \) and \( L/F \) a field extension satisfying \( [L : F] = r \). Let \( n = q^r - 1 \). Then \( t^n = 1 \) for all \( x \) in \( L \). The field \( L \) is a splitting field of \( t^n - 1 \) over \( F \), hence also of \( t^r - t \) in \( F[t] \), and \( G(K/F) \) is cyclic of order \( r \).
3. If \( F = \mathbb{Q} \), then \( [K : \mathbb{Q}] = \varphi(n) \) and \( G(K/\mathbb{Q}) \cong (\mathbb{Z}/n\mathbb{Z})^\times \).

We did not investigate the case when \( \text{char } F \) divides \( n \) in the above. Can you say anything in this case?

**Corollary 54.7.** Let \( F \) be a field of characteristic zero or \( \text{char } F \nmid n \) satisfying the polynomial \( t^n - 1 \) in \( F[t] \) splits over \( F \) with \( U \) the set of \( n \)th roots of units in \( F \). Suppose that \( K \) is a splitting field of the polynomial \( t^n - a \) in \( F[t] \) over \( F \) with \( a \) nonzero. Then \( G(K/F) \) is a cyclic subgroup isomorphic to a subgroup of \( U \). In particular, \( |G(K/F)| \mid n \).

**Proof.** Let \( U = \{z_1, \ldots, z_n\} \), a cyclic subgroup of \( F^\times \) of order \( n \). As \((t^n - a)' = nt^{n-1}\) has only zero as a root, \( t^n - a \) has \( n \) distinct roots in \( K \). Let \( r \) in \( K \) be a root of \( t^n - a \), then \( rz_i, i = 1, \ldots, n \), are all its roots. As \( z_i \in F \) for all \( i \), we know that \( K = F(r) \) and \( \sigma(rz_i) = \sigma(r)z_i \) for all \( \sigma \in G(K/F) \) and all \( i \). Consequently, \( \sigma(r) \) determines \( \sigma \) for each \( \sigma \in G(K/F) \). If \( \sigma(r) = rz_i \), then \( \sigma(r)/r = z_i \) lies in \( U \). Therefore, we have a map

\[
\varphi : G(K/F) \to U \quad \text{defined by} \quad \sigma \mapsto \frac{\sigma(r)}{r}.
\]

If \( \sigma \) and \( \tau \) lie in \( G(K/F) \), say \( \sigma(r) = rz_i \) and \( \tau(r) = rz_j \), then

\[
\varphi(\sigma \tau) = \frac{\sigma \tau(r)}{r} = \frac{\sigma(r)z_j}{r} = \frac{\sigma(r)z_j}{r} = \frac{\sigma(r)z_j}{r} = \frac{\sigma(r)}{r} = \frac{\tau(r)}{r} = \varphi(\sigma)\varphi(\tau),
\]

so \( \varphi \) is a well-defined group homomorphism. If

\[
\frac{\sigma(r)}{r} = \frac{\tau(r)}{r}, \quad \text{then} \quad \sigma(r) = \tau(r),
\]

hence \( \sigma = \tau \) and \( \varphi \) is monic. As \( U \) is cyclic of order \( n \), we have \( G(K/F) \) is cyclic of order dividing \( n \). \( \square \)

With these preliminaries, we can make some explicit computations.

**Computation 54.8.** Let \( F \) be a field. [You should fill in all the omitted details.]
1. $G(F/F) = 1$ and $F/F$ is Galois.

2. Suppose that the characteristic of $G$ is not two, $d$ an element in $F$ that is not a square in $F$. Then $K = F(\theta)$ is a splitting field of $t^2 - d$ in $F[t]$ over $F$ where $\theta$ is a root of $t^2 - d$. Then $K/F$ is Galois and $G(K/F) \cong \mathbb{Z}/2\mathbb{Z}$. [Of course, one often writes $\sqrt{d}$ for $\theta$.]

3. Let $f = (t^2 - 2)(t^2 - 3)$ in $Q[t]$ and $K = Q(\sqrt{2}, \sqrt{3})$. Then $K$ is a splitting field of $f$ over $Q$. As $\sqrt{2} \notin Q(\sqrt{3})$, the polynomial $t^2 - 2$ is irreducible in $Q(\sqrt{3})[t]$. Let $G(K/Q(\sqrt{3})) = \langle \sigma \rangle$ were $\sigma$ maps $\sqrt{3} \mapsto -\sqrt{3}$. Similarly, $G(K/Q(\sqrt{2})) = \langle \tau \rangle$ where $\tau$ maps $\sqrt{2} \mapsto -\sqrt{2}$. Then $G(K/Q) = \{1, \sigma, \tau, \sigma\tau\} \cong \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2\mathbb{Z}$ and $K/Q$ is Galois.

4. Let $f = t^3 - 2$ in $Q[t]$. It is irreducible, so $[Q(\sqrt[3]{2}) : Q] = 3$. Let $\zeta = \cos(2\pi/3) + \sqrt{-1}\sin(2\pi/3)$ and $K = Q(\sqrt[3]{2}, \zeta)$. As $m_Q(\zeta) = t^3 + t + 1$, we have $[Q(\zeta) : Q] = 2$. The field $K$ is a splitting field of $f$ over $Q$ so $[K : Q] = 6$. As $K$ is also the splitting field of $f$ over $Q(\zeta)$ and of degree three, the Galois group $G(K/Q(\zeta))$ is cyclic and isomorphic to a subgroup of $\mathbb{Z}/3\mathbb{Z}$. Since $f$ remains irreducible in $Q(\zeta)[t]$ (why?), there exists a $Q(\zeta)$-automorphism $\tau$ of $K$ satisfying $\sqrt[3]{2} \mapsto \zeta\sqrt[3]{2}$. Therefore, $G(K/Q(\zeta))$ is cyclic of order three. The polynomial $m_Q(\zeta)$ has no real roots, so remains irreducible in $Q(\sqrt[3]{2})[t]$, so $K$ is a splitting field of $t^3 + t + 1$ over $Q(\sqrt[3]{2})$ and there exists a $Q(\sqrt[3]{2})$-automorphism $\sigma : K \to K$ such that $\zeta \mapsto \zeta^{-1} = \overline{\zeta}$. Therefore, $G(K/Q(\sqrt[3]{2}))$ is cyclic of order two. It follows $G(K/Q)$ contains an element of order two and an element of order three, hence $[K : Q] = 6 = |G(K/Q)|$ and $K/Q$ is Galois. As

\[
\begin{align*}
\tau \sigma(\sqrt[3]{2}) &= \tau(\sqrt[3]{2}) = \zeta\sqrt[3]{2} \quad \text{and} \\
\sigma \tau(\sqrt[3]{2}) &= \sigma(\zeta\sqrt[3]{2}) = \zeta^{-1}\sqrt[3]{2},
\end{align*}
\]

$G(K/Q)$ is not abelian, hence isomorphic to $S_3 \cong D_3$. Any element of $G(K/Q)$ must take $\sqrt[3]{2}$ to a root of $t^3 - 2$ in $K$, so $Q(\sqrt[3]{2}) \subset \mathbb{R}$ means that the subgroup $G(Q(\sqrt[3]{2})/Q) = 1$ and $Q(\sqrt[3]{2})/Q$ is not Galois. Note also that $G(K/Q(\zeta))$ is normal in $G(K/Q)$ but $G(K/Q(\sqrt[3]{2}))$ is not normal in $G(K/Q)$. 

The following picture summarizes this:

\[
\begin{array}{c}
K = \mathbb{Q}(\sqrt[4]{2}, \zeta) \\
\mathbb{Q}(\sqrt[4]{2}) \quad \text{gal} \quad 2 \\
\text{not gal} \\
\mathbb{Q} \quad 3 \\
\mathbb{Q}(\zeta) \quad \text{gal} \quad 6
\end{array}
\]

5. Let \( f = t^5 - 1 \) in \( \mathbb{Q}[t] \) and \( \zeta = \cos(2\pi/5) + \sqrt{-1}\sin(2\pi/5) \) in \( \mathbb{C} \).
Then \( K = \mathbb{Q}(\zeta) \) is the splitting field of \( f \) and \( [K : \mathbb{Q}] = 4 = \varphi(5) \), as \( m_{\mathbb{Q}}(\zeta) = t^4 + t^3 + t^2 + t + 1 \). We know that group homomorphism \( G(K/\mathbb{Q}) \to (\mathbb{Z}/5\mathbb{Z})^\times \) given by \( \sigma \mapsto i \mod 5 \) if \( \sigma \in G(K/\mathbb{Q}) \) satisfies \( \sigma(\zeta) = \zeta^i \) is monic, so \( |G(K/\mathbb{Q})| = 1, 2, \text{ or } 4 \). As \( K \) is the splitting field of the irreducible polynomial \( m_{\mathbb{Q}}(\zeta) \) in \( \mathbb{Q}[t] \), the Galois group \( G(K/\mathbb{Q}) \) acts transitively on its roots, hence \( |G(K/\mathbb{Q})| = 4 \).
Therefore, \( K/\mathbb{Q} \) is Galois and \( G(K/\mathbb{Q}) \cong (\mathbb{Z}/5\mathbb{Z})^\times \).

Let \( \sigma_4 : K \to K \) be the \( \mathbb{Q} \)-automorphism determined by \( \zeta \mapsto \zeta^4 \).
Then \( \sigma_4 \) takes \( \zeta \) to \( \zeta^{-1} = \zeta^{-1} \), so has order two. If \( H = \langle \sigma_4 \rangle \), a subgroup of index two in \( G(K/\mathbb{Q}) \), then its fixed field is
\[
K^H = \{ a + b(\zeta^2 + \zeta^3) \mid a, b \in \mathbb{Q} \}
= \{ a + b(\zeta^2 + \zeta^2) \mid a, b \in \mathbb{Q} \} = \mathbb{Q}(\sqrt{5}).
\]
[Can you show this?] We have \( [K : K^H] = 2 = [K^H : \mathbb{Q}] \), so \( K/\mathbb{Q} \) is a square root tower (which leads to the construction of a regular pentagon).

6. Let \( K \) be a splitting field of the irreducible polynomial \( f = t^4 - 2 \) in \( \mathbb{Q}[t] \). The roots of \( f \) are: \( \sqrt[4]{2}, -\sqrt[4]{2}, -\sqrt{-1}\sqrt[4]{2}, -\sqrt{-1}\sqrt[4]{2}; \) so \( K = \mathbb{Q}(\sqrt[4]{2}, \sqrt{-1}) \). As \( \sqrt{-1} \notin \mathbb{Q}(\sqrt[4]{2}) \subset \mathbb{R} \), the polynomial \( t^4 + 1 \) is irreducible in \( \mathbb{Q}(\sqrt[4]{2})[t] \) and
\[
[K : \mathbb{Q}] = [K : \mathbb{Q}(\sqrt[4]{2})][\mathbb{Q}(\sqrt[4]{2}) : \mathbb{Q}] = 8 \geq |G(K/\mathbb{Q})|.
\]
Since \( t^4 - 1 \) splits over \( \mathbb{Q}(\sqrt{-1}) \), the group \( G(K/\mathbb{Q}(\sqrt{-1})) \) is cyclic of order 1, 2, or 4 by Corollary 54.7. We show that it is of order 4. We know that \( [K : \mathbb{Q}(\sqrt{-1})] = 4 \) and \( K \) is a splitting field of \( m_{\mathbb{Q}}(\sqrt[4]{2}) \).
It follows that \( m_{\mathbb{Q}}(\sqrt[4]{2}) \) must remain irreducible in \( \mathbb{Q}(\sqrt{-1})[t] \), i.e., \( m_{\mathbb{Q}(\sqrt{-1})}(\sqrt[4]{2}) = m_{\mathbb{Q}}(\sqrt[4]{2}) \). As \( G(K/\mathbb{Q}(\sqrt{-1})) \) acts transitively on the roots of \( m_{\mathbb{Q}(\sqrt{-1})}(\sqrt[4]{2}) \), we must have \( G(K/\mathbb{Q}(\sqrt{-1})) \) is a group of order at least 4, hence 4.
Now let
\[ \tau : K \to K \] be the element in \( G(K/Q(\sqrt[4]{2})) \)
satisfying \( \sqrt{-1} \mapsto -\sqrt{-1} \)
and
\[ \sigma : K \to K \] be the element in \( G(K/Q(\sqrt{-1})) \)
satisfying \( (\sqrt[4]{2}) \mapsto \sqrt{-1} \sqrt[4]{2} \).

Then \( \langle \tau, \sigma \rangle \subset G(K/Q)\). It is easy to see this forces
\[ |G(K/Q)| = 8. \]
Checking \( \tau \sigma \tau^{-1} = \sigma^3 \) shows that \( G(K/Q) \cong D_4 \). Note also that
\( K/Q \) is Galois but \( Q(\sqrt[4]{2})/Q \) is not as
\[ \langle Q(\sqrt[4]{2})/Q \rangle = \{1, \sigma^2\}_{Q(\sqrt[4]{2})}. \]
The following picture summarizes this:

\[
\begin{array}{ccc}
K & \xrightarrow{\text{gal}} & Q(\sqrt{-1}, \sqrt[4]{2}) \\
Q(\sqrt[4]{2}) & \xrightarrow{\text{not gal}} & \sqrt{-1} & \xrightarrow{\text{gal}} & Q(\sqrt{-1}) \\
\end{array}
\]

Note we have another interesting intermediate field of \( K/F \), viz.,
\( E = Q(\sqrt[4]{2}, \sqrt{-1}) \). As \([K : Q(\sqrt{-1})] = [K : Q]/[Q(\sqrt{-1}) : Q]\)
and \( \sqrt[4]{2} \) lies in \( K \) with \( \sqrt[4]{2} \) not in \( E \). \( Q < Q(\sqrt{-1}) < E < K \) is a square root tower. [Can you say more?]

7. Let \( p \) be a prime, \( F = (\mathbb{Z}/p\mathbb{Z})(x) \) with \( x \) transcendental over \( \mathbb{Z}/p\mathbb{Z} \)
and \( f = t^p - x \) a polynomial in \( F[t] \). Let \( K \) be a splitting field of \( f \)
over \( F \) and \( \alpha \) in \( K \) a root of \( f \). Then \( t^p - x = t^p - \alpha^p \) in \( K[t] \), so \( f \)
has only one root and \( K = F(\alpha) \) with \( \alpha \notin F \) (why?). Since \( G(K/F) \)
takes the roots of \( f \) to roots of \( f \), the Galois group \( G(K/F) = 1 \) and
\( K = K^{G(K/F)} \). So \( K/F \) is not Galois. Check that \( f \) is irreducible
over \( F \) so \([K : F] = p \) and \( K/F \) is a splitting field of an irreducible polynomial. Note that \( K/F \) is not separable.

Exercises 54.9.

1. Let \( K = Q(r) \) with \( r \) a root of \( t^3 + t^2 - 2t - 1 \in Q[t] \). Let \( r_1 = r^2 - 2 \).
Show that \( r_1 \) is also a root of this polynomial. Find \( G(K/Q) \).

2. Suppose the \(|K| = p^n \), \( p \) a prime, and \( F \subset K \). Show that \(|F| = p^m \)
for some \( m \) with \( m \mid n \). Moreover, \( G(K/F) \) is generated by the
Frobenius automorphism \( \alpha \mapsto \alpha^{p^m} \). In particular, \( G(K/F) \) is cyclic.
3. If $F$ is a finite field, $n$ a positive integer, then there exists an irreducible polynomial $f \in F[t]$ of degree $n$.

4. Let $F$ be a subfield of the real numbers, $f$ an irreducible quartic over $F$. Suppose that $f$ has exactly two real roots. Show that the Galois group of $f$ is either $S_4$ or of order 8.

5. Suppose that $K/F$ is Galois with Galois group $G(K/F) \cong S_n$. Show that $K$ is the splitting field of an irreducible polynomial in $F[t]$ of degree $n$ over $F$.

55. Galois Extensions

The object of Galois theory is to connect intermediate fields of a finite field extension to subgroups of the Galois group $G(K/F)$. The results and computations of the last section indicated a close relation when $K/F$ was a splitting field of a separable polynomial. We begin by pursuing that here.

**Definition 55.1.** We call a finite field extension $K/F$ normal if $K$ is the splitting field over $F$ of some non-constant polynomial in $F[t]$.

We shall see that the important properties of normal extensions are independent of a polynomial for which it is a splitting field except for whether the polynomial is separable or not.

It is useful to generalize the definition of a normal field extension. If $\{f_i\}_I$ is a set of non-constant polynomials in $F[t]$, we call an algebraic extension $K$ of $F$ a splitting field of the set of polynomials $\{f_i\}_I$, if every $f_i, i \in I$, splits in $K$ and $K$ is the smallest algebraic extension with this property. Such a splitting field exists and is unique up to an $F$-isomorphism. (Cf. Exercise 55.22(1)). For example, the splitting field of the set of all minimal polynomials over $F$ is algebraically closed and called an algebraic closure of $F$. If an algebraic extension $K$ of $F$ is the splitting field of a set of non-constant polynomials in $F[t]$, we say that $K/F$ is normal. Note that if $\{f_i\}_I$ is a finite set of non-constant polynomials in $F[t]$, then $K$ is a splitting field of this set of polynomials over $F$ if and only if it is the splitting field of $\prod_I f_i$.

The following is a very useful criterion to check if a finite field extension is normal.

**Proposition 55.2.** Let $K/F$ be a finite extension of fields. Then $K/F$ is normal if and only if any irreducible polynomial in $F[t]$ having a root in $K$ splits over $K$.

**Proof.** ($\Rightarrow$): Let $K$ be a splitting field the non-constant polynomial $g$ in $F[t]$ over $F$ and $f \in F[t]$ an irreducible polynomial having
a root $\alpha$ in $K$. Let $L/K$ with $\beta \in L$ a root of $f$. We must show that $\beta$ lies in $K$. As $f$ is irreducible, there exists an $F$-isomorphism $\tau : F(\alpha) \to F(\beta)$ satisfying $\alpha \mapsto \beta$. As $K(\alpha)$ is a splitting field of $g$ over $F(\alpha)$ and $K(\beta)$ is a splitting field of $g$ over $F(\beta)$, there exists an isomorphism $\sigma : K(\alpha) \to K(\beta)$ lifting $\tau$. In particular, $\sigma$ is also an $F$-isomorphism, so we have

$$[K(\alpha) : F] = [K(\alpha) : F(\alpha)][F(\alpha) : F] = [K(\beta) : F(\beta)][F(\beta) : F] = [K(\beta) : F].$$

But $\alpha \in K$, so $K(\alpha) = K$, and

$$[K : F] = [K(\alpha) : F] = [K(\beta) : F] = [K(\beta) : K][K : F].$$

Therefore, $[K(\beta) : K] = 1$, i.e., $\beta \in K$, as needed.

$\leftarrow\Rightarrow$: As $K/F$ is finite, there exist $\alpha_1, \ldots, \alpha_r$ in $K$, some $r$, with $K = F(\alpha_1, \ldots, \alpha_r)$ and each $\alpha_i$ algebraic over $F$. Clearly, $K$ is the splitting field of $f = \prod m_F(\alpha_i)$ over $F$, as each $m_F(\alpha_i)$ has a root in $K$, hence splits over $K$ and any splitting field of $f$ over $F$ in $K$ must contain all the $\alpha_i$'s.

Remark 55.3. 1. $\mathbb{Q}(\sqrt[3]{2})/\mathbb{Q}$ is not normal as the irreducible polynomial $t^3 - 2$ in $\mathbb{Q}[t]$ does not split over $\mathbb{Q}(\sqrt[3]{2})$

[Note: the general proof above is really the same as the special case done before.]

2. Let $p$ be a positive prime and $F = (\mathbb{Z}/p\mathbb{Z})(x)$ with $x$ transcendental over $\mathbb{Z}/p\mathbb{Z}$. Then $f = t^p - x$ in $F[t]$ is irreducible. (Cf. Exercise 52.10 (3).) Let $K/F$ with $\alpha \in K$ a root of $f$, i.e., $\alpha^p = x$. Set $E = F(\alpha)$. Then $f = (t - \alpha)^p$ in $E[t]$, hence $E$ is a splitting field of $f$ over $F$. Therefore, $E/F$ is normal. We saw that $G(E/F) = 1$, hence $E/F$ is not Galois.

Proposition 55.4. Let $K/F$ be a finite, normal extension and $K/E/F$ an intermediate field. Suppose that $\varphi : E \to K$ is an $F$-homomorphism. Then there exists an $F$-automorphism $\sigma \in G(K/F)$ satisfying $\sigma|_E = \varphi$, i.e., every $F$-homomorphism $E \to K$ arises from an $F$-automorphism of $K$ by restriction.

Proof. Let $K$ be a splitting field of the non-constant polynomial $f$ in $F[t]$, hence a splitting field of $f$ over $E$. In addition, $K$ is a splitting field of $f = \tilde{\varphi}(f)$ over $\varphi(E)$, so there exists an automorphism $\sigma : K \to K$ lifting $\varphi : E \to \varphi(E)$, so $\sigma$ lies in $G(K/F)$. \hfill \Box

Using Zorn’s Lemma, one can show that the last two propositions hold if we only assume that $K/F$ is algebraic. (Cf. Exercises 55.22(2),
(3), respectively.) Alternatively, one can use the existence of an algebraic closure (Corollary 50.2) and its lifting property (Theorem 50.3), together with a restriction argument giving the uniqueness of splitting fields.

**Definition 55.5.** Let $K/E_i/F$ be finite field extensions, $i = 1, 2$ with $K/F$ normal. We say that $E_1$ and $E_2$ are conjugate over $F$ if there exists an automorphism $\sigma \in G(K/F)$ satisfying $\sigma|_{E_1} : E_1 \to E_2$ is an isomorphism. Equivalently, by Proposition 55.4, there exists an $F$-isomorphism $E_1 \to E_2$.

**Remarks 55.6.** Let $L/F$ be a finite normal extension, $L/K/E_i/F$ field extensions, $i = 1, 2$, with $K/F$ also normal. Then $E_1$ and $E_2$ are conjugate relative to the extensions $K/E_i/F$ if and only if they are conjugate relative to the extensions $L/E_i/F$, as any $\hat{\sigma} \in G(L/F)$ lifts to some $\tilde{\sigma} \in G(L/F)$ and any $\tilde{\sigma} \in G(L/F)$ restricts to $\tilde{\sigma}|_K \in G(K/F)$. That $\hat{\sigma}$ is an $F$-automorphism of $K$ follows from the following:

**Corollary 55.7.** Let $K/F$ be a finite normal field extension and $K/E/F$ an intermediate field extension. Then $K/E$ is normal. Moreover, the following are equivalent:

1. $E/F$ is normal.
2. $\sigma(E) = E$ for every $\sigma \in G(K/F)$.
3. $\sigma|_E \in G(E/F)$ for every $\sigma \in G(K/F)$.
4. The map $\Phi : G(K/F) \to G(E/F)$ given by $\sigma \mapsto \sigma|_E$ is well-defined and an epimorphism.
5. $E$ is the only conjugate of $E$.

**Proof.** If $K$ is the splitting field of the non-constant polynomial $f \in F[t]$ over $F$, then it is also the splitting field of $f$ over $E$, so the first statement follows.

We turn to the equivalences. By Proposition 55.4, any $\varphi \in G(E/F)$ lifts to an element in $G(K/F)$ and if $G(K/F) \to G(E/F)$ by $\sigma \mapsto \sigma|_E$ is well-defined – usually it is not as, in general, $\sigma(E) \not\subseteq E$ – the equivalences of (2), (3), (4) and (5) become clear. So it suffices to show that (1) and (2) are equivalent.

(1) $\Rightarrow$ (2): Let $x \in E$ and $\sigma \in G(K/F)$. As $x \in E$ and $E/F$ is normal, we know that $m_F(x)$ splits over $E$. Since $G(K/F)$ takes the roots of $m_F(x)$ to roots of $m_F(x)$, the root $\sigma(x)$ also lies in $E$. As the argument also works for $\sigma^{-1}$, Statement (2) follows.

(2) $\Rightarrow$ (1): Let $f \in F[t]$ be an irreducible polynomial having a root $\alpha \in E$. As $K/F$ is normal, $f$ splits over $K$. Let $\beta \in K$ be root of $f$. We must show that $\beta \in E$, i.e., $f$ splits over $E$. We know that there exists
an $F$-isomorphism $\varphi : F(\alpha) \to F(\beta)$ taking $\alpha \mapsto \beta$. By Proposition 55.4, there exists an automorphism $\sigma \in G(K/F)$ satisfying $\varphi = \sigma|_{F(\alpha)}$, so $\beta = \varphi(\alpha) = \sigma(\alpha)$ lies in $E$ as needed. \hfill $\square$

Note in the above if $E/F$ is normal, then ker $\Phi = G(K/E) \triangleleft G(K/F)$ and $G(E/F) \cong G(K/F)/G(K/E)$.

**Definition 55.8.** Let $K/F$ be a finite field extension. Then a field extension $L$ of $K$ is called a normal closure of $K/F$ if $L/F$ is normal and $[L : K]$ is minimal with respect to this property. If this is the case, we write $L/K$ is a normal closure of $K/F$.

**Question 55.9.** How would you define the normal closure of an algebraic extension of $F$?

**Proposition 55.10.** Let $K/F$ be a finite field extension. Then a normal closure of $K/F$ exists and is unique up to a $K$-isomorphism.

**Proof.** Let $K = F(\alpha_1, \ldots, \alpha_n)$ with each $\alpha_i$ algebraic over $F$, $f = \prod m_F(\alpha_i)$ in $F[t]$, and $L$ a splitting field of $f$ over $K$, hence also a splitting field of $f$ over $F$. Therefore, $L/F$ is normal. By the uniqueness of splitting fields, $L$ is unique up to a $K$-isomorphism (relative to $f$). If $E/K$ is a normal extension, then each $m_F(\alpha_i)$ must split over $E$, so $f$ must split over $E$. The result now follows easily. \hfill $\square$

If in the above proposition, $K/F$ is also separable, then $L/F$ is also separable by Exercise 52.10(6). We shall also give a proof of this below.

The main technical result needed to establish the field theoretic equivalence of a Galois extension (rather than a group theoretic one) is the following:

**Proposition 55.11.** Let $K/F$ be a finite extension of fields of degree $n$. Suppose that $L/K$ with $L/F$ finite and normal and

$$\tau_1, \ldots, \tau_m : K \to L$$

are all the distinct $F$-homomorphisms. Then the following are true:

1. $m \leq n$.
2. $m = n$ if and only if $K/F$ is separable.
3. $\tau_1(K), \ldots, \tau_m(K)$ are all the conjugates of $K$ over $F$ in $L$ (and independent of the normal extension $L/F$).
4. Let $E = F(\tau_1(K) \cup \cdots \cup \tau_m(K))$. Then $E/K$ is a normal closure of $K/F$.

Using the exercises in 55.22, the assumption in the proposition that $L/K$ be finite may be dropped.
Proof. We may assume that \( \tau_1 \) is the inclusion, so \( m \geq 1 \). Let \( S = \{ \tau_1, \ldots, \tau_m \} \).

(1): By assumption, \( F \subset K^S \), so \( [K : F] \geq |S| = m \) by Artin’s Lemma 53.8.

(2): We induct on \( n \). The case of \( n = 1 \) is immediate, so we may assume that \( n > 1 \). Let \( \alpha \in K \setminus F \) be chosen with \( \alpha \) not separable, i.e., \( m_F(\alpha) \) not separable, if \( K/F \) is not separable. Let

\[
 r = [F(\alpha) : F] > 1
\]

and

\[
 \alpha = \alpha_1, \ldots, \alpha_{r_0} \text{ the distinct roots of } m_F(\alpha) \text{ in } L.
\]

[Note that \( m_F(\alpha) \) splits over \( L \) as \( L/F \) is normal.]

So we have

\[
 r = r_0 \text{ if and only if } \alpha \text{ is separable over } F
\]

if and only if \( K/F \) is separable.

In particular,

\[
 r_0 < r \text{ if } K/F \text{ is not separable.}
\]

Let

\[
 \varphi_i : F(\alpha) \to F(\alpha_i) [\subset L]
\]

be the \( F \)-isomorphism satisfying \( \alpha \mapsto \alpha_i, \ 1 \leq i \leq r_0 \).

Since the roots of \( m_F(\alpha) \) go to roots of \( m_F(\alpha) \), we have

\[
 \varphi_i, \ 1 \leq i \leq r_0, \text{ are all the } F\text{-homomorphisms } F(\alpha) \to L.
\]

As \( L/F \) is normal, each \( \varphi_i \) lifts to some

\[
 \hat{\varphi}_i \in G(L/F), \ 1 \leq i \leq r_0.
\]

[There can be many distinct lifts of each \( \varphi_i \) — chose one.]

Thus

\[
 \hat{\varphi}_i|_{F(\alpha)} \neq \hat{\varphi}_j|_{F(\alpha)} \text{ if } i \neq j.
\]

So we have the following:

(i) \( [K : F(\alpha)] = \frac{n}{r} \).

(ii) \( K/F(\alpha) \) is separable if \( K/F \) is separable.

By induction, we have:

There exist \( m_0 \) distinct \( F(\alpha) \)-homomorphisms

\[
 \psi_1, \ldots, \psi_{m_0} : K \to L
\]

with \( m_0 = \frac{n}{r} \) if and only if \( K/F(\alpha) \) is separable.

[Note if \( K/F \) is separable so is \( K/F(\alpha) \).]
Let 
\[ \rho_{ij} = \widehat{\varphi}_i \psi_j : K \rightarrow L, \ 1 \leq i \leq r_0, \ 1 \leq j \leq m_0. \]
and
\[ S = \{ \rho_{ij} \mid 1 \leq i \leq r_0, \ 1 \leq j \leq m_0 \}. \]
Each \( \rho_{ij} \) is an \( F \)-homomorphism, so \( S \subset \{ \tau_1, \ldots, \tau_m \} \).

Claim 55.12. \( S \) contains all the \( F \)-homomorphisms \( K \rightarrow L \) and \( m = |S| = r_0m_0 \).

If we establish the claim, then \( m = r_0m_0 \leq m_0r \leq n. \)

So \( K/F \) is separable if and only if \( r = r_0 \) and \( m_0 = \frac{n}{r} \) if and only if \( n = m_0r_0 = m. \)

In particular, if we prove the claim, (2) follows.

Let \( \rho : K \rightarrow L \) be an \( F \)-homomorphism. Then \( \rho(\alpha) \) is a root of \( m_F(\alpha) \), so there exists an \( i \) satisfying \( \rho|_{F(\alpha)} = \varphi_i \), hence \( \widehat{\varphi}_i^{-1}\rho : K \rightarrow L \) is an \( F(\alpha) \)-homomorphism. Consequently, there exist a \( j \) satisfying \( \psi_j = \widehat{\varphi}_i^{-1}\rho \). Therefore, \( \rho = \widehat{\varphi}_i \psi_j = \rho_{ij} \) lies in \( S \).

To finish proving the claim, we must show that the \( \rho_{ij} \)'s are distinct.

Suppose that \( \rho_{ij} = \rho_{kl} \). Then
\[ \varphi_i = \rho_{ij}|_{F(\alpha)} = \rho_{kl}|_{F(\alpha)} = \varphi_k, \]
hence \( \widehat{\varphi}_i = \widehat{\varphi}_k \), so \( i = k. \) Since
\[ \widehat{\varphi}_i \psi_j = \rho_{ij} = \rho_{kl} = \widehat{\varphi}_k \psi_l, \]
it follows that \( \psi_j = \psi_l \) as \( \widehat{\varphi}_i \) is an automorphism. Consequently, \( j = l. \)
This shows that the \( \rho_{ij} \)'s are distinct and establishes the claim.

(3). Each of the \( \tau_i : K \rightarrow L \) above lifts to some \( \widehat{\tau}_i \) in \( G(L/F) \) as before.
In particular, \( \widehat{\tau}_i(K) = \tau_i(K) \). If \( \sigma \in G(L/F) \), there exists a \( j \) such that \( \sigma|_K = \tau_j \). It follows that \( \tau_1(K), \ldots, \tau_m(K) \) are all the conjugates of \( K \) (repetitions possible).

(4): Suppose that \( L/F \) is also a normal closure of \( K/F \). Set \( E = F(\tau_1(K) \cup \cdots \cup \tau_m(K)) \subset L \). Write \( K = F(\beta_1, \ldots, \beta_s) \), for some \( \beta_i \), and set \( f = \prod_{i=1}^s m_F(\beta_i) \) in \( F[t] \). Then \( L \) is a splitting field of \( f \) over \( F \). By the argument proving (2), we know that for any root of \( m_F(\beta_i) \), there exist an \( F \)-homomorphism \( K \rightarrow L \) sending \( \beta_i \) to that root. As \( \tau_1, \ldots, \tau_m : K \rightarrow L \) are all the \( F \)-homomorphisms and each must take roots of \( m_K(\beta_i) \) to roots of \( m_F(\beta_i) \), the field \( E \) contains all the roots of the \( m_F(\beta_i) \). As \( L/F \) is a splitting field of \( f \), we must have \( E = L. \) \( \square \)
Corollary 55.13. Let $K/F$ be a extension of fields with $\alpha$ an element in $K$ algebraic over $F$. Then $\alpha$ is separable over $F$ if and only if $F(\alpha)/F$ is separable.

Proof. Let $L/F(\alpha)$ be a field extension such that $L/F$ is finite and normal. We know that $m_F(\alpha)$ splits over $L$. Let $\alpha = \alpha_1, \ldots, \alpha_m$ be the distinct roots of $m_F(\alpha)$ in $L$. We know that there exist $F$-isomorphisms $\varphi_i : F(\alpha) \to F(\alpha_i)$ satisfying $\alpha \mapsto \alpha_i$ for $i = 1, \ldots, m$ and these are all the distinct $F$-homomorphisms into $L$. Thus $F(\alpha)/F$ is separable if and only if $[F(\alpha) : F] = \deg m_F(\alpha) = m$ if and only if $\alpha$ is separable over $F$.

Proposition 55.14. Let $K$ be a splitting field of a separable polynomial in $F[t]$. Then $K/F$ is separable and Galois.

Proof. Suppose that $K$ is a splitting field of the separable polynomial $f \in F[t]$. We may assume that $F \subset K$. Let $\alpha \in K$ be a root of $f$ not in $F$. As $\alpha$ is separable over $F$, the extension $F(\alpha)/F$ is separable by the corollary. Therefore, there exist $m = [F(\alpha) : F]$ $F$-homomorphisms $\varphi_i : F(\alpha) \to K$, $1 \leq i \leq m$. As $K/F$ is normal, each $\varphi_i$ lifts to some $\widehat{\varphi}_i \in G(K/F)$. Since $K$ is a splitting field of the separable polynomial $f$ over $F(\alpha)$, by induction on $[K : F]$, we conclude that $K/F(\alpha)$ is separable. As $K/K$ is the normal closure of $K/F(\alpha)$ (as well as $K/F$), there exist precisely $[K : F(\alpha)]$ distinct $F(\alpha)$-homomorphisms $\psi_j : K \to K$, $1 \leq j \leq [K : F(\alpha)]$. Moreover, the argument to prove Claim 55.12 shows that $$S = \{\widehat{\varphi}_i \psi_j \mid 1 \leq i \leq [F(\alpha) : F], \ 1 \leq j \leq [K : F(\alpha)]\}$$ is the set of all $F$-homomorphisms $K \to K$ and has $[K : F] = [K : F(\alpha)][F(\alpha) : F]$ elements, hence $K/F$ is separable. Since $S \subset G(K/F)$, we have $|G(K/F)| = [K : F]$ by Artin’s Lemma 53.8, so $K/F$ is also Galois.

Corollary 55.15. Let $K/F$ be a finite separable field extension and $L/K$ a normal closure of $K/F$. Then $L/F$ is separable and Galois.

Proof. Let $K = F(\alpha_1, \ldots, \alpha_n)$. Then $L$ is a splitting field of the separable polynomial $\prod m_F(\alpha_i)$ in $F[t]$, hence $L/F$ is separable and Galois.

Corollary 55.16. Let $K = F(\alpha_1, \ldots, \alpha_n)$ with $\alpha_i$ separable over $F$ for each $i = 1, \ldots, n$. Then $K/F$ is separable.

Proof. Let $L/K$ be a normal closure of $K/F$, hence the splitting field of the separable polynomial $\prod m_F(\alpha_i)$. Therefore, $L/F$ is separable hence so is $K/F$. 

□
Corollary 55.17. If $K/F$ is finite field extension and $K/E/F$, then $K/F$ is separable if and only if $K/E$ and $E/F$ are separable.

**Proof.** ($\Rightarrow$) has already been done.
($\Leftarrow$): Let $L/K$ be a (finite) normal extension and 

$\varphi_1, \ldots, \varphi_m : E \to L$ be all the distinct $F$-homomorphisms 

$\psi_1, \ldots, \psi_n : K \to L$ be all the distinct $E$-homomorphisms. 

Let $\hat{\varphi}_i \in G(L/F)$ lift $\varphi_i$ for $1 \leq i \leq m$. Then as in the proof of Claim 55.12,

$$|\{\hat{\varphi}_i \psi_j \mid 1 \leq i \leq m, \ 1 \leq j \leq n\}| = mn.$$ 

As $K/E$ and $E/F$ are separable, $[K : F] = [K : E][E : F] = mn$, hence $K/F$ is separable by Proposition 55.11.

We arrive at our goal of finding the field theoretic interpretation for a finite field extension $K/F$ to be Galois.

**Theorem 55.18.** Let $K/F$ be a finite extension of fields. Then $K/F$ is Galois if and only if $K/F$ is normal and separable.

**Proof.** ($\Leftarrow$) has already been done.
($\Rightarrow$): Let $L/K$ be a normal closure of $K/F$. As $K/F$ is Galois, the number of $F$-homomorphisms $K \to L$ is at least $|G(K/F)| = [K : F]$. It follows by Proposition 55.11 that there exist exactly $|G(K/F)| = [K : F]$ such $F$-homomorphisms and, in addition, $K/F$ must be separable. Since each of these homomorphisms lies in $G(K/F)$, Proposition 55.11(4), shows that $L = F(\cup_{G(K/F)} \sigma(K)) = K$.

**Corollary 55.19.** Let $K/F$ be a finite extension and $K/E/F$ an intermediate field. Suppose that $K/F$ is Galois. Then 

1. $K/E$ is Galois.
2. $E/F$ is Galois if and only if $E/F$ is normal.

**Proof.** As $K/F$ is Galois, it is separable and normal, so $K/E$ and $E/F$ are separable and $K/E$ is normal. Therefore, $K/E$ is Galois and $E/F$ is normal if and only if $E/F$ is Galois.

Let $K/F$ be a finite field extension and $K/E/F$ an intermediate field. Recall that we have seen, in general, $K/F$ being Galois does not imply that $E/F$ is Galois. For example, $K = \mathbb{Q}(\sqrt[3]{2}, e^{2\pi i/3})$, $E = \mathbb{Q}(\sqrt[3]{2})$, and $F = \mathbb{Q}$. A question left to the reader is: If $K/E$ and $E/F$ are Galois, is $K/F$ Galois?

**Definition 55.20.** A field $F$ is called **perfect** if any algebraic extension $K/F$ is separable.
So we have

**Remark 55.21.** Let $F$ be a field.

1. If char $F = 0$, then $F$ is perfect.
2. If $F$ is a finite field, then $F$ is perfect.
3. If $F$ is perfect, $K/F$ a finite extension, then $K/F$ is Galois if and only if $K/F$ is normal.
4. If $x$ is transcendental over $\mathbb{Z}/p\mathbb{Z}$, with $p > 0$ a prime and $F = (\mathbb{Z}/p\mathbb{Z})(x)$, then $F$ is not perfect as $t^p - x$ in $F[t]$ is not separable.

**Exercises 55.22.**

1. Let $F$ be a field and $\{f_i\}_I$ be a set of non-constant polynomials in $F[t]$. Prove that a splitting field of $\{f_i\}_I$ over $F$ exists and is unique up to an $F$-isomorphism.

2. Let $K/F$ be an algebraic extension of fields. Then $K/F$ is normal if and only if any irreducible polynomial in $F[t]$ having a root in $K$ splits over $K$.

3. Let $K/F$ be a normal (possibly infinite) extension of fields and $K/E/F$ an intermediate field. Suppose that $\varphi : E \to K$ is an $F$-homomorphism. Then there exists an $F$-automorphism $\sigma \in G(K/F)$ satisfying $\sigma|_E = \varphi$, i.e., every $F$-homomorphism $E \to K$ arises from an $F$-automorphism of $K$ by restriction.

4. Let $K/F$ be an algebraic field extension. Answer Question 55.9 and prove that a normal closure of $K/F$ exists and is unique up to a $K$-isomorphism.

5. Let $F$ be a field and $L$ and algebraic closure of $F$. Show that $L/F$ is normal. Is it true that $L^{G(L/F)} = F$, i.e., $L/F$ is (infinite) Galois? If not, give an example when it is not.

6. Show that there exist uncountably many field automorphisms of $\mathbb{C}$.

7. Let $K/F$ be an algebraic extension of fields and $K/E/F$ an intermediate field. Show that $K/F$ is separable if and only if $K/E$ and $E/F$ are separable.

8. Suppose that $K/F$ is Galois and $\alpha \in K$ has precisely $r$ distinct images under $G(K/F)$. Show $[F(\alpha) : F] = r$.

9. Let $F$ be a field of positive characteristic $p$ and $f = t^p - t - a \in F[t]$.
   (a) Show the polynomial $f$ has no multiple roots.
   (b) If $\alpha$ is a root of $f$, show so is $\alpha + k$ for all $0 \leq k \leq p - 1$.
   (c) Show $f$ is irreducible if and only if $f$ has no root in $F$. 

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(d) Suppose that \( a \neq b^p - b \) for any \( b \in F \). Find \( G(K/F) \) where \( K \)

is a splitting field of \( t^p - t - a \in F[t] \).

10. Let \( K/F \) be a finite field extension and \( K/E/F \) an intermediate

field. Suppose that \( K/E \) and \( E/F \) are Galois. Is \( K/F \) Galois?

Either prove or provide a counterexample.

11. Prove if \( F \) is not a finite field and \( u, v \) are algebraic and separable

over \( F \) that there exists an element \( a \in F \) such that \( F(u + av) \). Is this true if \( |K| < \infty \) with \( K(u) < K(u_v) \) and \( K(v) < K(u_v) \)?

12. Let \( K = \mathbb{Q}(\sqrt{2}, \sqrt{3}, u) \) where \( u^2 = (9 - 5\sqrt{3})(2 - \sqrt{2}) \). Show that

\( K/Q \) is normal and find \( G(K/Q) \).

13. Let \( L/F \) be a finite normal extension and \( E = L^{G(L/F)} \). Show

\[ E = \{ \alpha \mid \alpha \text{ is purely inseparable over } F \} \]

(Cf. Exercises 52.10(11) and (14).) In particular, \( L/E \) is separable and \( E/F \) is purely inseparable, i.e., we have

\[
\begin{array}{c}
L \\
\downarrow \text{pinsep} \\
L_{\text{sep}} \\
\downarrow \text{sep} \\
F \\
\downarrow \text{pinsep}
\end{array}
\]

\[ E = L^{G(L/F)} \]

56. The Fundamental Theorem of Galois Theory

In this section we put together the pieces that we have developed,

i.e., the group theoretic and field theoretic interpretations of Galois

extensions. As an application we prove the Fundamental Theorem of

Algebra. We begin with the following lemma.

**Lemma 56.1.** Let \( K/F \) be a finite Galois extension of fields, \( H_i \subset \) \( G(K/F) \) a subgroup and \( E_i = K^{H_i} \) for \( i = 1, 2 \). If \( \sigma \) is an element of \( \) the Galois group \( G(K/F) \), then

\[ \sigma|_{E_i} : E_1 \to E_2 \text{ is an isomorphism if and only if } H_2 = \sigma H_1 \sigma^{-1}; \]

i.e., \( E_1 \) and \( E_2 \) are conjugate over \( F \) (via \( \sigma \)) if and only if \( H_1 \) and \( H_2 \) are conjugate in \( G(K/F) \) (via \( \sigma \)).

**Proof.** We know that \( K/E_i \) is Galois and \( H_i = G(K/E_i) \) for \( i = 1, 2 \). In particular, \( [K: E_i] = |H_i| \) for \( i = 1, 2 \), so under either condition, we have

\[ |H_1| = |H_2| = [K : E_1] = [K : E_2]. \]
(⇒): Let $\tau \in H_1$ and $y \in E_2$. Then $\sigma^{-1}(y)$ lies in $E_1$, so $\sigma\tau\sigma^{-1}(y) = \sigma\tau^{-1}(y) = y$ as $\tau|_{E_1} = 1_{E_1}$. Therefore, we have $\sigma H_1 \sigma^{-1} \subset H_2 = G(K/E_2)$, hence $\sigma H_1 \sigma^{-1} = H_2$ as $|\sigma H_1 \sigma^{-1}| = |H_1| = |H_2| < \infty$.

(⇐): Let $x \in E_1$ and $\tau \in H_1$. Then $\sigma \tau \sigma^{-1}(\sigma(x)) = \sigma \tau(x) = \sigma(x)$, as $\tau|_{E_1} = 1|_{E_1}$. Consequently, $\sigma(x)$ lies in $K^\sigma H_1 \sigma^{-1} = K H_2 = E_2$, and $\sigma|_{E_1} : E_1 \to E_2$ and is a homomorphism. As $[K : E_1] = [K : E_2]$ and $[E_1 : F] = [E_2 : F]$, the map $\sigma_{E_1}$ is an isomorphism (since a linear monomorphism of finite dimensional vector spaces of the same dimension).

\[\square\]

**Notation 56.2.** Let $K/F$ be an algebraic extension of fields. We set

\[
\mathcal{F}(K/F) := \{ E \mid K/E/F \text{ is an intermediate field} \}
\]

\[
\mathcal{G}(K/F) := \{ H \mid H \subset G(K/F) \text{ is a subgroup} \}.
\]

Putting together all our results leads to:

**Theorem 56.3.** (The Fundamental Theorem of Galois Theory) Suppose that $K/F$ is a finite Galois extension. Then

\[
i : \mathcal{G}(K/F) \to \mathcal{F}(K/F) \text{ given by } H \mapsto K^H
\]

is a bijection of sets with inverse

\[
j : \mathcal{F}(K/F) \to \mathcal{G}(K/F) \text{ given by } E \mapsto G(K/E).
\]

Moreover,

1. The bijection $i$ is order-reversing, i.e., $H_1 \subset H_2$ if and only if $K^H_1 \supset K^H_2$.
2. If $E \in \mathcal{F}(K/F)$, then $E/F$ is normal (hence Galois) if and only if $G(K/E) \triangleleft G(K/F)$.
3. If $E \in \mathcal{F}(K/F)$ with $E/F$ normal, then the canonical epimorphism

\[\sigma : G(K/F) \to G(E/F) \text{ given by } \sigma \mapsto \sigma|_E\]

induces an isomorphism

\[G(K/F)/G(K/E) \to G(E/F) \text{ given by } \sigma G(K/E) \mapsto \sigma|_E.\]

4. If $H \in \mathcal{G}(K/F)$, then $|H| = [K : K^H]$ and $K/K^H$ is Galois with $H = G(K/K^H)$.
5. If $H \in \mathcal{G}(K/F)$, then $[G(K/F) : H] = [K^H : F]$.

We have the following picture:
with the identified verticals having the same degree (index) and if $E/F$ is normal, $G(K/F)/G(K/E) \cong G(E/F)$.

**Proof.** We already know that $H_1 = H_2$ in $G(K/F)$ if and only if $K^{H_1} = K^{H_2}$ in $F(K/F)$, so $i : G(K/F) \to F(K/F)$ is injective. Let $E \in F(K/F)$. Then $K/E$ is Galois, so $G(K/E) \mapsto K^{G(K/E)} = E$. Hence $i : G(K/F) \to F(K/F)$ is bijective.

(1) is clear.

(2): We know that if $E \in F(K/F)$, then $E/F$ is normal if and only if $\sigma(E) = E$ for all $\sigma \in G(K/F)$. By Lemma 56.1, this is true if and only if $\sigma G(K/E)\sigma^{-1} = G(K/E)$ for all $\sigma \in G(K/F)$, i.e., if and only if $G(K/E) \triangleleft G(K/F)$.

We have seen if $E/F$ is normal then $\Phi : G(K/F) \to G(E/F)$ by $\sigma \mapsto \sigma_{|E}$ is a well-defined group epimorphism. As

$$\ker \Phi = \{\sigma \in G(K/F) \mid \sigma_{|E} = 1_E\} = G(K/E),$$

$\Phi$ induces an isomorphism $G(K/F)/G(K/E) \cong G(E/F)$.

(4): As $K/K^H$ is Galois, $H = G(K/K^H)$, so $H = [K : K^H]$ by Artin’s Theorem 53.15.

(5): By Artin’s Theorem 53.15,

$$[K : K^H][K^H : F] = [K : F] = |G(K/F)| = [G(K/F) : H]|H|,$$

so (4) $\Rightarrow$ (5). $\square$

**Corollary 56.4.** Let $K/F$ be a finite Galois extension of fields. Then $F(K/F)$ is finite, i.e., there exist only finitely many intermediate fields $E$ with $K/E/F$.

**Proof.** By the Fundamental Theorem of Galois Theory, we have $|F(K/F)| = |G(K/F)|$, hence $F(K/F)$ is a finite set. $\square$

**Corollary 56.5.** Let $K/F$ be a finite separable extension. Then $F(K/F)$ is finite.
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Proof. Let \( L/K \) be a normal closure of \( K/F \). Then \( L/F \) is finite, separable, and normal, so Galois. Thus \( |\mathcal{F}(K/F)| \leq |\mathcal{F}(L/F)| \) is finite.

We give a few examples, before deriving significant consequences of this fundamental theorem.

Example 56.6. 1. We know that \( K = \mathbb{Q}(\sqrt{2}, \sqrt{3}) \) is the splitting field of \((t^2 - 2)(t^2 - 3)\) in \( \mathbb{C} \) over \( \mathbb{Q} \). Then we have seen that \( G(K/\mathbb{Q}) = \{1, \sigma, \tau, \sigma\tau\} \cong \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2\mathbb{Z} \) with

\[
\sigma : K \to K \text{ the } \mathbb{Q}(\sqrt{3})\text{-automorphism given by } \sqrt{2} \mapsto -\sqrt{2},
\]

\[
\tau : K \to K \text{ the } \mathbb{Q}(\sqrt{2})\text{-automorphism given by } \sqrt{3} \mapsto -\sqrt{3}.
\]

Intermediate fields and corresponding subgroups are given by the following diagrams:

```
\begin{tikzpicture}
  \node (1) at (0,0) {1};
  \node (sigma tau) at (1,1) {\langle \sigma\tau \rangle};
  \node (sigma) at (1,-1) {\langle \sigma \rangle};
  \node (G) at (-1,0) {G(K/\mathbb{Q})};
  \node (K) at (0,3) {K};
  \node (Q2) at (-1,1) {\mathbb{Q}(\sqrt{2})};
  \node (sqrt6) at (-1,2) {\mathbb{Q}(\sqrt{6})};
  \node (sqrt3) at (-1,3) {\mathbb{Q}(\sqrt{3})};
  \node (Q) at (-1,0) {\mathbb{Q}};
  \draw (K) -- (Q2);
  \draw (K) -- (sqrt6);
  \draw (K) -- (sqrt3);
  \draw (Q2) -- (Q);
  \draw (sqrt6) -- (Q);
  \draw (sqrt3) -- (Q);
\end{tikzpicture}
```

2. Let \( K \) be a splitting field of \( t^4 - 2 \) over \( \mathbb{Q} \). Then \( K = \mathbb{Q}(\sqrt[4]{2}, \sqrt{-1}) \) and \( |K : \mathbb{Q}| = 8 \). If

\[
\sigma : K \to K \text{ the } \mathbb{Q}(\sqrt[4]{2})\text{-automorphism given by } \sqrt{-1} \mapsto -\sqrt{-1},
\]

\[
\tau : K \to K \text{ the } \mathbb{Q}(\sqrt{-1})\text{-automorphism given by } \sqrt[4]{2} \mapsto \sqrt{-1}\sqrt[4]{2},
\]

then \( |\langle \sigma \rangle| = 2, \ |\langle \tau \rangle| = 4, \ \sigma\tau\sigma^{-1} = \tau^3, \ G(K/\mathbb{Q}) = \langle \sigma, \tau \rangle \cong D_4 \). Intermediate fields and corresponding subgroups are given by the following diagrams with \( \theta = \sqrt[4]{2}, \ i = \sqrt{-1} \):
3. Let $L = F(t_1, \ldots, t_n)$, the quotient field of $F[t_1, \ldots, t_n]$. The symmetric group $S_n$ acts on $L$ as $F$-automorphisms by permuting the $t_i$, i.e.,

$$\sigma(t_i) = t_{\sigma(i)}$$

for all permutations $\sigma \in S_n$ and $i = 1, \ldots, n$.

We view $S_n \subset G(L/F)$. Let

$$s_j = s_j(t_1, \ldots, t_n) := \sum_{1 \leq i_1 < \cdots < i_j \leq n} t_{i_1} \cdots t_{i_j}$$

be the $j$th elementary symmetric polynomial of $t_1, \ldots, t_n$ for $j = 1, \ldots, n$ and set $s_0 = 1$. We know that $t_1, \ldots, t_n$ are the roots of the polynomial

$$f := t^n - s_1 t^{n-1} + \cdots + (-1)^n s_n$$

in $F[s_1, \ldots, s_n][t]$, as $f = (t - t_1) \cdots (t - t_n)$ in $L[t]$. Let $K = L^{S_n}$, so $L/K$ is Galois. Let $E := F(s_1, \ldots, s_n) \subset K$. Clearly, $L$ is a splitting field of $f$ over $E$, and by our theory, we know that

$$[L : E] \leq (\deg f)! = n!$$

and

$$[L : K] = |G(L/K)| = |S_n| = n!.$$

So \( K = E = F(s_1, \ldots, s_n) \) and \( G(L/K) = G(L/F(s_1, \ldots, s_n)) \cong S_n \). Note that \( f \) must be irreducible in \( E[t] \), lest \( f = f_1f_2 \) in \( E[t] \) with \( 0 < \deg f_1 < n \) for \( i = 1, 2 \), and we would have \( [L : E] \leq (\deg f_1)(\deg f_2)! < (\deg f)! = [L : E] \). Note further that \( s_1, \ldots, s_n \) are algebraically independent over \( F \) as \( t_1, \ldots, t_n \) are and \( K/E \) is finite. (Cf. Proposition 48.4.)

As an application of this computation, we also get a weak form of the Fundamental Theorem of Symmetric Polynomials. (Cf. Theorem ?? below.) Let \( F \) be a field and \( S_n \) act on \( F[t_1, \ldots, t_n] \) by restriction of the above action. Then we have with \( s_j \) the \( j \)th elementary symmetric function in \( t_1, \ldots, t_n \):

**Corollary 56.7.** Let \( F \) be a field and \( f \) be a symmetric polynomial in \( F[t_1, \ldots, t_n] \). Then \( f \) is a rational function in the elementary symmetric functions in \( t_1, \ldots, t_n \), i.e., \( f \in F(s_1, \ldots, s_n) \).

The Fundamental Theorem of Symmetric Functions [cf. Theorem 70.4 below] implies that a symmetric polynomial in \( F[t_1, \ldots, t_n] \) actually lies in \( F[s_1, \ldots, s_n] \), even with \( F \) replaced by a commutative ring. [Cf. Theorem 70.4 for the proof.]

We continue with the Examples 56.6

4. Every finite group \( G \) is isomorphic to some Galois group, i.e., there exists a Galois field extension \( L/K \) with \( G \cong G(L/K) \):

Let \( G \) be a finite group of order \( n \) and \( F \) be an arbitrary field. Set \( L = F(t_1, \ldots, t_n) \). If \( E = F(s_1, \ldots, s_n) \) as in previous example. Then \( L/E \) is Galois with \( G(L/E) \cong S_n \). By Cayley’s Theorem 12.5, there exists a group monomorphism \( \lambda : G \to S_n \), so \( G \cong \lambda(G) \subset S_n \). Let \( K = L^{\lambda(G)} \), then \( L/K \) is Galois with \( G(L/K) \cong G \).

**Open Problem 56.8.** (Inverse Galois Problem) Let \( H \) be a finite group. Does there exist a Galois extension \( L \) of \( Q \) with \( G(L/Q) \cong H ? \)

The answer to this question is known for some types of groups, e.g., cyclic groups, abelian groups, solvable groups, \( S_n, A_n \).

We now turn to some significant consequences of the Fundamental Theorem of Galois Theory.

**Theorem 56.9.** (Primitive Element Theorem) Let \( K/F \) be a finite extension of fields. If \( F(K/F) \) is finite, then there exists an element \( \alpha \) in \( K \) such that \( K = F(\alpha) \). In particular, this is true if \( K/F \) is separable.

**Proof.** Case 1. \( F \) is finite:
$K$ must also be a finite field, so $K^\times = \langle \alpha \rangle$ for some $\alpha \in K$. Therefore, $K = F(\alpha)$ and $\mathcal{F}(K/F)$ is finite.

**Case 2.** $F$ is infinite:

Choose $\alpha \in K$ with $[F(\alpha) : F]$ maximal, so $[F(\alpha) : F] \leq [K : F]$. Suppose that $F(\alpha) < K$, then there exists a $\beta \in K \setminus F(\alpha)$. For each $a \in F$ consider the subfield $F(\alpha + a\beta) < F(\alpha, \beta) \subset K$. As $\mathcal{F}(F(\alpha, \beta)/F)$ is a subset of the finite set $\mathcal{F}(K/F)$ and $F$ is infinite, there exists elements $c$ and $d$ in $F$ with $c \neq d$ satisfying $F(\alpha + c\beta) = F(\alpha + d\beta)$. Since $c - d$ lies in $F^\times$ and $(c - d)\beta$ lies in $F(\alpha + c\beta)$, we must have $\beta \in F(\alpha + c\beta)$ and hence that $\alpha$ also lies in $F(\alpha + c\beta)$. It follows that $F(\alpha, \beta) \subset F(\alpha + c\beta)$, contradicting the maximality of $F(\alpha)$. □

By Exercise 52.10(9b), we know this is not always true if $K/F$ is not separable.

**Remark 56.10.** Let $F$ be an infinite field and $K/F$ a finite extension with $\mathcal{F}(K/F)$ finite. Suppose that $\alpha$ and $\beta$ lie in $K$ and $S$ is an infinite subset of $F$. Then the proof above shows that there exists an element $c$ in $S$ such that $F(\alpha, \beta) = F(\alpha + c\beta)$. In particular, if $K = F(\alpha_1, \ldots, \alpha_n)$ and $S$ is an infinite subset of $F$, then there exist infinitely many $n$-tuples $c_1, \ldots, c_n$ in $S^n$ satisfying $K = F(c_1\alpha_1 + \cdots + c_n\alpha_n)$.

We finally prove the loose end in the constructibility of regular $n$-gons by straight-edge and compass (cf. Theorem 51.8).

**Theorem 56.11.** (Square Tower Theorem) Let $F$ be a field of characteristic different from two and $K/F$ a finite normal field extension. Then $K/F$ is a square root tower if and only if $[K : F] = 2^e$ for some positive integer $e$.

**Proof.** ($\Rightarrow$) has previously been established.

($\Leftarrow$): Suppose that $[K : F] = 2^e$, for some integer $e$. Let $K = F(\alpha_1, \ldots, \alpha_n)$. We know that $\deg m_F(\alpha_i) = [F(\alpha_i) : F] = [K : F]$ for each $i$. As $\text{char } F \neq 2$, each $m_F(\alpha_i)$ is separable, hence the extension $K/F$ is separable. Consequently, $K/F$ is normal and separable, hence Galois, so $|G(K/F)| = [K : F] = 2^e$. We know that there exists a subgroup $H$ of the 2-group $G(K/F)$ of index two, therefore, $H$ is a normal subgroup. It follows that $K^H/F$ is normal of degree two. By Remark 51.7(6), $K^H/F$ is a square root tower, as $\text{char } F$ is not two. The field extension $K/K^H$ is Galois of degree $2^{e-1}$, so a square root tower by induction on $e$. Therefore, $K/F$ is a square root tower by Remark 51.7(4). □

Next we give a proof of the Fundamental Theorem of Algebra that uses a minimal amount of analysis. Indeed the only analysis that we
need is the Intermediate Value Theorem from Calculus (which is really the completeness of the real line).

**Theorem 56.12.** (Fundamental Theorem of Algebra) *The field of complex numbers is algebraically closed.*

**Proof. Claim.** $\mathbb{C} = \mathbb{C}^2 := \{x^2 \mid x \in \mathbb{C}\}$:

Let $a$ be a positive real number. The polynomial $t^2 - a$ in $\mathbb{R}[t]$ has a real root by the Intermediate Value Theorem, since $f(1 + a) > 0$ and $f(a) < 0$. Therefore, every positive real number is a square in $\mathbb{R}$ hence in $\mathbb{C}$. Let $-a = (\sqrt{-1})^2 a$, it follows that every real number is a square in $\mathbb{C}$. Let $\alpha = a + b\sqrt{-1}$ with $a$ and $b$ real numbers. We must show $\alpha$ is a square in $\mathbb{C}$. We may assume that $b \neq 0$. Let $x = (a + \sqrt{a^2 + b^2})/2$ in $\mathbb{R}$. Then $x$ is positive (why?), so there exists a real number $c$ such that $x = c^2$. Similarly, $y = (-a + \sqrt{a^2 + b^2})/2$ in $\mathbb{R}$ is positive, so there exists a real number $d$ such that $y = d^2$. Then $\alpha = (c + d\sqrt{-1})^2$, establishing the claim.

[The above arises from $\sqrt{\alpha} = \sqrt{|\alpha|}e^{\sqrt{-1}\theta/2}$ with $\theta$ the angle satisfying $\tan \theta = b/a$ and the half angle formula in trigonometry – which we do not need!]

Now let $f$ be a non-constant polynomial in $\mathbb{C}[t]$. We must show that $f$ splits over $\mathbb{C}$. Let $K/\mathbb{C}$ be a splitting field of $f$ and $L/K$ a normal closure of $K/\mathbb{R}$. If we show that $L = \mathbb{C}$, we are done. We know that $L/\mathbb{R}$ is Galois and $2 \mid [\mathbb{C} : \mathbb{R}]$. Therefore, there exists a Sylow 2-subgroup $H$ of $G(K/\mathbb{R})$. In particular, $[L^H : \mathbb{R}] = [G(L/\mathbb{R}) : H]$ is odd. Since $\mathbb{R}$ is a field of characteristic zero, it is perfect, so there exists an element $\alpha \in L^H$ satisfying $L^H = \mathbb{R}(\alpha)$ by the Primitive Element Theorem 56.9. The element $\alpha$ in $L^H$ is a root of the irreducible polynomial $m_{\mathbb{R}}(\alpha)$ in $\mathbb{R}[t]$ of odd degree. By the Intermediate Value Theorem, any real polynomial of odd degree has a real root. Therefore, $\alpha$ is real, so $L^H = \mathbb{R}(\alpha) = \mathbb{R}$. Therefore, $G(L/\mathbb{R}) = G(L/L^H)$ is a 2-group. As $L/\mathbb{R}$ is Galois and $\text{char} \ \mathbb{R} \neq 2$, $L/\mathbb{R}$ is a square root tower by the Square Root Tower Theorem. It follows that $L(\mathbb{C})/\mathbb{C}$ is also a square root tower by Remark 51.7(2). By the claim, $\mathbb{C}$ has no proper square root towers over it, so $L(\mathbb{C}) = \mathbb{C}$ as needed. \[\square\]

**Exercises 56.13.**

1. Show that the lattices of subgroups and subfields in Example 56.6(2) are correct.

2. Let $F$ be a perfect field and $K/F$ a finite Galois extension. Show if $G(K/F) \cong S_n$, then $K$ is the splitting field of an irreducible polynomial of degree $n$. Is this still true if $G(K/F) \cong A_n$?
3. Let $K$ be a splitting field of $f \in \mathbb{Q}[t]$. Find $K$, $G(K/\mathbb{Q})$, and all intermediate fields if:
   (i) $f = t^4 - t^2 - 6$.
   (ii) $f = t^3 - 3$.
4. Let $K$ be a splitting field of $t^5 - 2 \in \mathbb{Q}[t]$.
   (i) Find $G(K/\mathbb{Q})$.
   (ii) Show that there exists a group monomorphism $G(K/\mathbb{Q}) \to S_5$.
   (iii) Find all subgroups of $G(K/\mathbb{Q})$ and the corresponding fields.
5. Suppose that $L/F$ is a finite Galois extension and $L/K/F$ an intermediate field. Show that $G(K/F) = N_{G(L/F)}(G(L/K))/G(L/K))$, where $N_{G(L/F)}(G(L/K))$ is the normalizer of $G(L/K)$ in $G(L/F)$.
6. Suppose that $K/F$ is a finite Galois extension of fields. Let $F \subset E \subset K$ and $L$ the smallest subfield of $K$ containing $E$ such that $L/F$ is normal. Show $G(K/L) = \bigcap_{\sigma \in G(K/E)} \sigma G(K/E) \sigma^{-1}$, the core of $G(K/E)$ in $G(K/F)$.
7. Let $L/F$ be a finite Galois extension, $L/K_i/F$ intermediate fields, and $H_i = G(K_i/F)$ for $i = 1, 2$. Show $H_1 \cap H_2 = G(L/K_1(K_2))$ and the fixed field of the smallest group in $G(L/F)$ containing $H_1$ and $H_2$ is $K_1(K_2)$.
8. Let $K/F$ be a finite Galois extension. Suppose that $p^r \nmid [K : F]$ but $p^{r+1} \nmid [K : F]$. Show that there exist fields $L_i$, $1 \leq i \leq r$, satisfying $F \subset L_r < L_{r-1} < \ldots < L_1 < L_0 = K$ with $L_i/L_{i+1}$ is normal, $[L_i : L_{i+1}] = p$ and $p \not| [L_r : F]$ for each $i$.
9. Let $f$ be an irreducible quartic over a field $F$ of characteristic zero, $G$ the Galois group of $f$ (i.e., the Galois group of $K/F$ with $K$ a splitting field of $f$ over $F$), $u$ a root of $f$. Show that there is no field properly between $F$ and $F(u)$ if and only if $G = A_4$ or $G = S_4$ and that there exist such irreducible polynomials with Galois group $A_4$ and $S_4$. [This will explode the myth that there must be an intermediate field when the dimension is not prime.]
10. Let $L/F$ and $K/F$ be finite extensions of fields with $L/F$ Galois and $K$, $L$ lying in some extension field $F$. Let $LK = K(L) = L(K)$. Show that $LK/K$ and $L/L \cap K$ are Galois and the restriction map $\varphi : G(LK/K) \to G(L/F)$ given by $\sigma \mapsto \sigma|_L$ a monomorphism with image $G(L/L \cap K)$. In particular, if $L \cap K = F$, then $G(LK/K) \cong G(L/F)$.
11. Suppose that $M/F$ be a finite Galois extension and $K$ and $L$ two intermediate fields of $M/F$ with $K/F$ Galois. Let $KL = L(K)$, $N = G(M/K) \triangleleft G(M/F)$, and $H = G(M/L)$. Then $KL/L$ is
Galois by Exercise 10. Show, in addition, that all of the following are true:
(a) \( G(M/KL) = H \cap N \triangleleft H \).
(b) \( G(KL/L) \cong H/H \cap N \cong HN/N \).
(c) \( G(K/K \cap L) \cong HN/N \).

12. Let \( E/F \) be an extension of fields and \( E/K_i/F \) be intermediate fields, \( i = 1, 2 \) with \( K_i/F \) finite Galois for \( i = 1, 2 \). Let \( K_1K_2/F \) be a Galois extension and the restriction map \( \varphi : G(K_1K_2/F) \to G(K_1/F) \times G(K_2/F) \) given by \( \sigma \mapsto (\sigma|_{K_1}, \sigma|_{K_2}) \) is a monomorphism. In particular, it is an isomorphism if \( K_1 \cap K_2 = F \).

13. Let \( L/F \) be a finite Galois extension. Suppose that \( G(L/F) = G_1 \times \cdots \times G_n \), a direct product of subgroups. Let \( K_i \) be the fixed field of \( G_1 \times \cdots \times \{1\} \times \cdots \times G_n \) where the group \( \{1\} \) occurs in the \( i \)th place. Show that \( K_i/F \) is Galois for every \( i \), \( K_i \cap F \cup \{K_j \mid j \neq i\} = F \), and \( L = F(\cup K_i) \).


15. Let \( F \) be a field having no nontrivial field extensions of odd degree and \( K/F \) a finite field extension. Show if \( K \) has no field extensions of degree two, then \( F \) is perfect and \( K \) is algebraically closed. (Cf. Exercise 52.10(11)).

57. Addendum: Infinite Galois Theory

In this section, we indicate how Galois Theory for finite extensions extends to arbitrary algebraic extensions. In particular, normal extensions need not be finite and an algebraic extension \( K \) of \( F \) is called Galois if \( F = K^{G(K/F)} \). We must generalize some of our proofs.

**Proposition 57.1.** Let \( K/F \) be an algebraic extension. If \( K/F \) is normal, and \( \sigma : K \to K \) an \( F \)-homomorphism, then \( \sigma \) lies in \( G(K/F) \), i.e., \( \sigma \) is onto.

**Proof.** Let \( \alpha \) lie in \( K \). Then the minimal polynomial \( m_F(\alpha) \) splits over \( K \) (just as in the finite normal case by Exercise 55.22(2)), so there exists an intermediate field \( K/E/F \) with \( E \) a splitting field of \( m_F(\alpha) \) over \( F \). Thus \( E/F \) is a finite normal extension, so \( \sigma(E) = E \), i.e., \( \sigma|_E \in G(E/F) \). Consequently, there exists a \( \beta \) in \( E \), hence \( K \), such that \( \sigma(\beta) = \alpha \). \( \square \)

**Proposition 57.2.** Let \( K/F \) be an algebraic extension. If \( K/F \) is normal and \( K/E/F \) is an intermediate field, then any \( F \)-homomorphism \( \sigma : E \to K \) lifts to an element of \( G(K/F) \).
Theorem 57.3. Let $K/F$ be an algebraic extension. If $K/F$ is normal, then $K/F$ is Galois if and only if $K/F$ is separable and normal.

Proof. $(\Leftarrow)$: Let $\alpha$ lie in $K^{G(K/F)}$. As above there exists an intermediate field $K/E/F$ with $E$ a splitting field of $m_F(\alpha)$ over $F$. As $K/F$ is separable, so is $E/F$, hence $E/F$ is finite normal and separable so Galois. Let $\sigma$ lie in $G(E/F)$, By the last proposition, there exists an extension $\widehat{\sigma}$ of $\sigma$ in $G(K/F)$, i.e., $\widehat{\sigma}|_E = \sigma$, so by assumption $\sigma(\alpha) = \widehat{\sigma}(\alpha) = \alpha$. Therefore, $\alpha$ lies in $E^{G(E/F)} = F$.

$(\Rightarrow)$: Let $\alpha$ lie in $F$. Every element of $S = \{\sigma(\alpha) \mid \sigma \in G(K/F)\}$ is a root of $m_F(\alpha)$ (which a priori may have other and/or multiple roots), so $S$ is a finite set. Let $\alpha = \alpha_1, \alpha_2, \ldots, \alpha_n$ be the finite distinct elements of $S$ and $f = \prod_{i=1}^{n}(t - \alpha_i)$. As $G(K/F)$ permutes $\alpha_1, \ldots, \alpha_n$, we must have $f = \widehat{\sigma}f$ for every $\sigma \in G(K/F)$, i.e., $f$ lies in $F[t]$, But then $f \mid m_f(\alpha)$ in $F[t]$, so $f = m_F(\alpha)$ is separable and $F = m_F(\alpha)$ splits over $K$. It follows that $K/F$ is separable and normal.

Remark 57.4. Let $K/F$ be an algebraic Galois extension of fields, $H$ a subgroup of $G(K/F)$, and $E = K^H$. Suppose that $\sigma$ and $\tau$ are elements of $G(K/F)$. If $\tau^{-1}\sigma$ lies in $H$, then for every $x \in E$, we have $\tau^{-1}\sigma(x) = x$, so $\sigma|_E = \tau|_E$. However, if we only know that $(\tau^{-1}\sigma)|_E = 1_E$, we can only conclude that $\tau^{-1}\sigma$ lies in $G(K/F)$. It does not necessarily follow that $\tau^{-1}\sigma$ lies in $H$. Therefore, we only know that $H \subset G(K/E) = G(K/K^H)$. In the case that $K/F$ is a finite extension, we get equality by counting using Artin’s Theorem, but this does not work in the arbitrary algebraic case.

We do, however, have the following:

Lemma 57.5. Let $K/F$ be an algebraic and Galois extension of fields and $K/E/F$ an intermediate field. If $[E : F]$ finite, then the index $[G(K/F) : G(K/E)]$ is finite and $[G(K/F) : G(K/E)] = [E : F]$.

Proof. Suppose that $\sigma_1, \ldots, \sigma_n : E \to K$ are all the distinct $F$-homomorphisms. These lift to $\widehat{\sigma}_1, \ldots, \widehat{\sigma}_n$ in $G(K/F)$, i.e., $\widehat{\sigma}_i|_E = \sigma_i$ for each $i$, by Proposition 57.2 and, as $E/F$ is finite separable, we have $n = [E : F]$ by Proposition 55.11. Suppose that $\widehat{\sigma}$ is an element of $G(K/F)$. Then there exists an $i$ such that $\widehat{\sigma}|_E = \sigma_i$. Therefore, $(\widehat{\sigma}_i)^{-1}\widehat{\sigma}$ lies in $G(K/E)$. It follows that

\begin{equation}
G(K/F) = \bigvee_{i=1}^{n} \widehat{\sigma}_i G(K/E),
\end{equation}
and the result follows.

We now state without proof the major new result needed. The ideas involve topological knowledge that we have not assumed. We summarize some of the concepts that we need. A topology \( \mathcal{T} \) on a space \( X \) is a collection of sets, called open sets closed under arbitrary unions and finite intersections. \( X \) is then called a topological space (via \( \mathcal{T} \)). Let \( X \) have a topology \( \mathcal{T} \). A closed set in \( X \) is one that is the complement of an open set in \( \mathcal{T} \). If \( x \) is a point in \( X \), an open set containing \( x \) is called a neighborhood of \( x \) (e.g., an open interval containing \( x \) if \( x \) lies in \( X = \mathbb{R} \)). A subcollection of \( \mathcal{T} \) consisting of neighborhoods of \( x \) such that every neighborhood of \( x \) contains an element of this subcollection is called a fundamental system of neighborhoods of \( x \) (e.g., all the open intervals with rational end points containing \( x \) if \( X = \mathbb{R} \)). A map \( f : X \to Y \) of topological spaces is called continuous if the preimage of an open set in \( Y \) under \( f \) is open in \( X \). Then the following is true:

**Theorem 57.7.** (Krull) Suppose that \( K/F \) is an algebraic Galois extension of fields. Let

\[
\mathcal{N}(K/F) := \{ G(K/E) \mid K/E/F \text{ with } E/F \text{ finite Galois} \}.
\]

Then there exists a topology on \( G \) compatible with the group structure of \( G \) (i.e., the group operation on \( G \) is continuous relative to this topology) and which has \( \mathcal{N}(K/F) \) as a fundamental system of neighborhoods of the identity \( 1_K \) in \( G(K/F) \). With this topology, \( G(K/F) \) is a Hausdorff, totally disconnected, topological group.

Let \( X \) be a topological space. It is called Hausdorff if given any two distinct points in \( X \), there exist disjoint neighborhoods of these two points and compact if every collection of open sets that cover \( X \), i.e., whose union is \( X \), contains a finite subcollection that covers \( X \). Finally, \( X \) is called totally disconnected if there exists a subcollection of the topology consisting of sets that are both open and closed and such that every open set can be obtained from this subcollection by taking arbitrary unions of finite intersections from it. The topology in the theorem is called a profinite topology.

In the terminology of the theorem, if \( \sigma \) is an element of \( G(K/F) \), then by translation a fundamental system of neighborhoods for \( \sigma \) is just \( \{ \sigma H \mid H \in \mathcal{N}(K/F) \} \). By equation 57.6, if \( H \in \mathcal{N}(K/F) \), then \( [G(K/F) : H] \) is finite and every element of \( \mathcal{N}(K/F) \) is both open and closed (look at complements). This is the totally disconnected statement. If \( H \) is a subgroup of \( G(K/F) \), let \( \overline{H} \) denote the closure of \( H \) in \( G(K/F) \) in this topology, i.e., the smallest closed set in \( G(K/F) \) containing \( H \). This means that if \( \sigma \) in \( G(K/F) \) lies in \( \overline{H} \), then every
open neighborhood of $\sigma$ intersects $H$ non-trivially. Because of the topology on $G(K/F)$, this means that $H \cap \sigma H'$ is nonempty for every $H'$ in $\mathcal{N}(K/F)$.

**Lemma 57.8.** Let $K/F$ be an algebraic Galois extension of fields, $H$ a subgroup of $G(K/F)$, and $E = K^H$. Then $G(K/E) = \overline{H}$, the closure of $H$ in $G(K/F)$.

**Proof.** $G(K/E) \subset \overline{H}$: Let $\sigma$ lie in $G(K/E)$. We must show that $\sigma \in \overline{H}$. As remarked above, since $\sigma \mathcal{N}(K/F)$ is a fundamental system of neighborhoods for $\sigma$, it suffices to show that $H \cap \sigma H'$ is nonempty for all $H'$ in $\mathcal{N}(K/F)$. Suppose that $H' \in \mathcal{N}(K/F)$. As $K^{H'}/F$ is finite Galois, it is finite separable, so $K^{H'} = F(\alpha)$ for some $\alpha$ in $K$ by the Primitive Element Theorem 56.9. Let $K/L/E$ be an intermediate field such that $L/E$ is finite, normal, and separable, hence finite Galois. Set $H' = G(K/L)$. By assumption, $H \cap \sigma H'$ is nonempty, so there exists an element $\tau$ in $H \cap \sigma H'$. In particular, as $\alpha$ lies in $L$, we have $\sigma(\alpha) = \tau(\alpha)$, so $\sigma^{-1}\tau$ lies in $H'$. This shows that $\tau \in \sigma H' \cap H$.

$\overline{H} \subset G(K/E)$: Let $\sigma$ lie in $\overline{H}$. We must show that $\sigma$ lies in $G(K/E)$, i.e., if $\alpha$ is an element of $E$, then $\sigma(\alpha) = \alpha$. Let $K/L/F(\alpha)$ with $L/F$ finite, normal, and separable, hence finite Galois. Set $H' = G(K/L)$. By assumption, $H \cap \sigma H'$ is nonempty, so there exists an element $\tau$ in $H \cap \sigma H'$. In particular, $\tau|_E = 1_E$, since $E = K^H$, and $\sigma(\alpha) = \tau(\alpha)$, since $F(\alpha) \subset K^{H'}$. Therefore, $\sigma(\alpha) = \tau(\alpha) = \alpha$. \qed

The Fundamental Theorem of Galois Theory in this more general setting becomes:

**Theorem 57.9.** Let $K/F$ be an algebraic Galois extension of fields and

$$G_c(K/F) := \{H | H \text{ is a closed subgroup of } G(K/F)\},$$

then

$$i : G_c(K/F) \to \mathcal{F}(K/F) \text{ given by } H \mapsto K^H$$

is an order reversing bijection.

**Proof.** By the lemma, we know that

$$G_c(K/F) = \{G(K/E) | E \in \mathcal{F}(K/F)\}.$$
i is injective: If \( K^{H_1} = K^{H_2} \), then by the lemma,
\[
H_1 = \overline{H_1} = G(K/K^{H_1}) = G(K/H_2) = \overline{H_2} = H_2.
\]

i is surjective: Let \( E \) lie in \( \mathcal{F}(K/F) \), so \( E \subset K^{G(K/E)} \). We must show that \( G(K/E) \) moves \( K \setminus E \). Let \( \alpha \in K \setminus E \). As before, let \( K/L/E \) with \( L/E \) finite normal, hence Galois, and \( \alpha \in L \). As \( \alpha \) does not lie in \( E \), there exist an element \( \sigma \) in \( G(L/E) \) satisfying \( \sigma(\alpha) \neq \alpha \). By Proposition 57.2, there exists a lift \( \widehat{\sigma} \) in \( G(K/F) \) of \( \sigma \).

\[ \square \]

Remark 57.10. Let \( K/F \) be an algebraic Galois extension of fields, \( E \in \mathcal{F}(K/F) \). Then

\( i \) \( E/F \) is normal if and only if \( G(K/E) \triangleleft G(K/F) \).

\( ii \) If \( E/F \) is normal, then \( G(E/F) \cong G(K/F)/G(K/E) \).

We leave the proofs of these as exercises.

Example 57.11. Let \( S = \{ \sqrt{a} \mid a > 0 \text{ a prime or } a = -1 \} \subset \mathbb{C} \) and \( K \) the splitting field of \( \{ t + a^2 \mid a \in S \} \) in \( \mathbb{C} \). Then \( K/\mathbb{Q} \) is Galois, and it is easy to see that \( K \) is the compositum of all quadratic extensions of \( \mathbb{Q} \) in \( \mathbb{C} \), i.e., \( K = \mathbb{Q}(\bigcup_{d \in \mathbb{Z}} \mathbb{Q}(\sqrt{d})) \), the number of which is countable. [More generally, if \( L/F \) is a field extension, \( L/K_i/F, i \in I \), intermediate extensions, then \( F(\bigcup_i K_i) \) is called the compositum of the \( K_i, i \in I \), in \( L \).] If \( \sigma \) and \( \tau \) lie in \( G(K/\mathbb{Q}) \), then we have \( \sigma^2 = 1_K \) and \( \sigma \tau = \tau \sigma \), i.e., \( G(K/\mathbb{Q}) \) is abelian and all non-identity elements are of order two. It follows that \( G(K/\mathbb{Q}) \) is a vector space over \( \mathbb{Z}/2\mathbb{Z} \), hence has a \( \mathbb{Z}/2\mathbb{Z} \)-basis, say \( \mathcal{B} \).

Let \( T \) be a subset of \( S \) and define \( \sigma : S \to S \) by
\[
\sigma(x) = \begin{cases} 
  x & \text{if } x \in T \\
  -x & \text{if } x \notin T.
\end{cases}
\]

Then \( \sigma \) induces an element in \( G(K/\mathbb{Q}) \). In particular, \( G(K/\mathbb{Q}) \) is uncountable, so cannot have a countable \( \mathbb{Z}/2\mathbb{Z} \)-basis. Therefore, for all but countably many \( x \in \mathcal{B} \), the subgroup \( H_x := \langle \mathcal{B} \cup \{ x \} \rangle \) is a normal subgroup of index two in \( G(K/\mathbb{Q}) \) that does not correspond to a quadratic extension of \( \mathbb{Q} \), i.e., only countable many \( H_x \) are closed (or open) in \( G(K/\mathbb{Q}) \).

Exercises 57.12.

1. Let \( K/F \) be Galois and
\[
\mathcal{N} := \{ G(K/E) \mid K/E/F \text{ with } E/F \text{ finite Galois} \}.
\]
Show that \( \bigcap_{\mathcal{N}} G(K/E) = 1 \).
2. Assuming Theorem 57.7, show that the topology is Hausdorff.
3. Show the two items in Remark 57.10 are true.

58. Roots of Unity

An arithmetic function $f : \mathbb{Z}^+ \to \mathbb{C}$ is called a multiplicative function if $f(1)$ is nonzero and $f(mn) = f(m)f(n)$, whenever $m$ and $n$ are relatively prime. For example, the Euler phi-function is multiplicative (by the Chinese Remainder Theorem). Note if $f$ is a multiplicative function, then $f(1) = 1$. Let $I : \mathbb{Z}^+ \to \mathbb{C}$ be the multiplicative function defined by $I(1) = 1$ and $I(n) = 0$ for all $n > 1$. The set $\mathcal{A} := \{ f \mid f : \mathbb{Z}^+ \to \mathbb{C} \text{ with } f(1) \neq 0 \}$ is an abelian monoid with unity $I$ under the Dirichlet product $\star : \mathcal{A} \times \mathcal{A} \to \mathbb{C}$ defined by

$$(f \star g)(n) := \sum_{d|n} f(d)g\left(\frac{n}{d}\right)$$

(cf. Exercise 58.21(1)). This group arises from the study of Dirichlet series, i.e., absolutely convergent series of the form $\sum_{n} \frac{f(n)}{n^s}$, $f$ a multiplicative function and $s$ a complex variable. It can be shown that $\mathcal{M} := \{ f \mid f \in \mathcal{A} \text{ multiplicative} \}$ is a subgroup. Basic to this is Möbius inversion that we now investigate.

Definition 58.1. The Möbius function $\mu : \mathbb{Z}^+ \to \mathbb{C}$ is defined by $\mu(1) = 1$ and if $n = p_1^{e_1} \cdots p_r^{e_r}$ is a standard factorization of $n > 1$, then

$$\mu(n) := \begin{cases} (-1)^r & \text{if } n \text{ is square-free, i.e., } e_1 = \cdots = e_r = 1 \\ 0 & \text{otherwise} \end{cases}$$

Let $U : \mathbb{Z}^+ \to \mathbb{C}$ be the multiplicative function defined by $U(n) = 1$ for all $n$.

Lemma 58.2. The Möbius function $\mu$ is multiplicative and satisfies:

1. $I = \mu \star U$ i.e.,

$$\sum_{d|n} \mu(d) = \begin{cases} 1 & \text{if } n = 1 \\ 0 & \text{if } n > 1. \end{cases}$$

2. If $n = p_1^{e_1} \cdots p_r^{e_r}$ is a standard factorization of $n > 1$, then

$$\sum_{d|n} |\mu(d)| = 2^r.$$

Proof. Let $m$ and $n$ be relatively prime positive integers. Clearly, either $m$ or $n$ is not square free if and only if $mn$ is not square free. It follows easily that $\mu$ is multiplicative. Let $\epsilon = \mu$ or $|\mu|$. If $n > 1$
is not square-free, then \( \epsilon(n) = 0 \). In particular, we may assume that \( n = p_1 \cdots p_r \) is a standard factorization of \( n > 1 \). Then

\[
\sum_{d|n} \epsilon(d) = \epsilon(1) + \sum_i \epsilon(p_i) + \sum_{i_1 < i_2} \epsilon(p_i_1 p_i_2) + \\
\cdots + \sum_{i_1 < \cdots < i_j} \epsilon(p_i_1 \cdots p_i_j) + \cdots + \epsilon(p_1 \cdots p_r)
\]

\[
= 1 + \epsilon(p_1) \binom{r}{1} + \epsilon(p_1 p_2) \binom{r}{2} + \cdots + \epsilon(p_1 \cdots p_r) \binom{r}{r}.
\]

If \( \epsilon = \mu \), this is \( 1 - \binom{r}{1} + \binom{r}{2} + \cdots + (-1)^r \binom{r}{r} = (1 - 1)^r = 0 \) and if \( \epsilon = |\mu| \), this is \( 1 + \binom{r}{1} + \binom{r}{2} + \cdots + \binom{r}{r} = (1 + 1)^r = 2^r \). \( \square \)

**Proposition 58.3.** (Möbius Inversion Formula) Let \( f \) and \( g \) be arithmetic functions not zero at 1. Then

\[
f(n) = \sum_{d|n} g(d) \text{ if and only if } g(n) = \sum_{d|n} f(d) \mu\left(\frac{n}{d}\right).
\]

**Proof.** The proposition is equivalent to \( f = g \ast U \) if and only if \( g = f \ast \mu \). By Exercise 8.5(15), the Dirichlet product is associative with \( I \) a unity, so \( f = g \ast U \) if and only if \( f \ast \mu = (g \ast U) \ast \mu = g \ast (U \ast \mu) = g \ast I = g \). \( \square \)

For each \( n \), let \( \zeta_n \) be a fixed primitive \( n \)th root of unity in \( \mathbb{C} \). Define the \( n \)th **cyclotomic polynomial** \( \Phi_n(t) \) in \( \mathbb{C}[t] \) by \( \Phi_n = \prod_{\zeta_n} (t - \zeta) \) in \( \mathbb{C}[t] \), a polynomial of degree \( \varphi(n) \). Clearly,

\[
t^n - 1 = \prod_{d|n} \Phi_d,
\]

so \( \log(t^n - 1) = \sum_{d|n} \log \Phi_d \). By the Möbius Inversion Formula (viewing the cyclotomic polynomials as polynomial functions on the integers or use Exercise 58.21(6)), we see that

\[
\Phi_n = \prod_{d|n} (t^d - 1)^{\mu\left(\frac{n}{d}\right)} = \prod_{d|n} (t^{\frac{n}{d}} - 1)^{\mu(d)} \text{ lies in } \mathbb{C}[t] \cap \mathbb{Q}(t) = \mathbb{Q}[t]
\]

and is monic. By induction \( \Phi_d \) lies in \( \mathbb{Z}[t] \) for \( d < n \), hence so does \( \Phi(n) \) by Corollary 34.6 (the corollary to Gauss’s Lemma).
Example 58.4. Using the formula above, we have
\[ \Phi_{12} = \prod_{d \mid 12} (t^d - 1)^{\mu(d)} \]
\[ = (t^{12} - 1)^{\mu(1)}(t^6 - 1)^{\mu(2)}(t^4 - 1)^{\mu(3)}(t^3 - 1)^{\mu(4)} \]
\[ = (t^2 - 1)^{\mu(6)}(t - 1)^{\mu(1)} \]
\[ = (t^{12} - 1)(t^2 - 1) = (t^6 - 1)(t^4 - 1) = t^4 - t^2 + 1. \]

Definition 58.5. Let \( K/F \) be a finite Galois extension of fields. We say that \( K/F \) is an abelian (respectively, cyclic) extension if \( G(K/F) \) is abelian (respectively, cyclic).

Theorem 58.6. Let \( \zeta \) be a primitive \( n \) root of unity in \( \mathbb{C} \). Then \( \mathbb{Q}(\zeta)/\mathbb{Q} \) is an abelian extension of degree \( \phi(n) \).

Proof. We must show that \( \Phi_n \) is irreducible in \( \mathbb{Q}[t] \), equivalently in \( \mathbb{Z}[t] \). Let \( \zeta \) be a primitive \( n \)th root of unity. As \( \mathbb{Z}[t] \) is a UFD, we see that there exist polynomials \( f \) and \( g \) in \( \mathbb{Z}[t] \) satisfying \( \Phi_n = fg \) in \( \mathbb{Z}[t] \) with \( f \) irreducible in \( \mathbb{Z}[t] \) and having \( \zeta \) as a root. As \( \Phi_n \) is monic, we may assume that both \( f \) and \( g \) are monic. [Note that this implies that \( f = m\mathbb{Q}(\zeta) \).] Let \( p \) be a prime such that \( p \not| n \).

Claim: \( \zeta^p \) is a root of \( f \).

If we prove the claim, it would follow that \( \zeta^m \), with \( m \) relatively prime to \( n \), is also root of \( f \), hence \( f = \Phi_n \) establishing the theorem.

Suppose the claim is false. We can write \( t^n - 1 = fh \) in \( \mathbb{Z}[t] \) with \( h \in \mathbb{Z}[t] \) monic by Corollary 34.6 (the corollary to Gauss’s Lemma). As \( \zeta^p \) is not a root of \( f \), it must be a root of \( h \). This implies that \( \zeta \) is a root of the monic polynomial \( h_p := h(t^p) \) in \( \mathbb{Z}[t] \). It follows that \( h_p = fg \) in \( \mathbb{Z}[t] \) for some monic \( g \) in \( \mathbb{Z}[t] \) by Corollary 34.6. Let \( \overline{\cdot} : \mathbb{Z} \to \mathbb{Z}/p\mathbb{Z} \) be the canonical epimorphism; and, if \( f \) lies in \( \mathbb{Z}[t] \), let \( \overline{f} \) be the image of \( f \) in \( (\mathbb{Z}/p\mathbb{Z})[t] \). Applying the Frobenius homomorphism shows that
\[ \overline{f^p} = \overline{f}(t^p) = \overline{h_p} = \overline{f} \overline{g} \]
and has \( \overline{\zeta} \) as a root, hence so does \( \overline{h_p} \). Therefore, \( t^n - 1 = \overline{f} \overline{h_p} \) has a multiple root over \( \mathbb{Z}/p\mathbb{Z} \). As \( p \not| n \), this is impossible. This establishes the claim and the theorem. \[ \square \]

Remark 58.7. Let \( p \) be an odd prime and \( \zeta \) a primitive \( p^r \)th root of unity. Then \( G(\mathbb{Q}(\zeta)/\mathbb{Q}) \cong (\mathbb{Z}/p^r\mathbb{Z})^\times \) is cyclic of order \( p^{r-1}(p-1) \). In particular, \( \mathbb{Q}(\zeta) \) contains a unique quadratic extension of \( \mathbb{Q} \) using Proposition E.4.
Corollary 58.8. Let \( \zeta_m, \zeta_n, \) and \( \zeta_{mn} \) be primitive \( m \text{th}, n \text{th}, \) and \( (mn) \text{th} \) roots of unity over \( \mathbb{Q} \) with \( m \) and \( n \) relatively prime positive integers. Then \( \mathbb{Q}(\zeta_m, \zeta_n) = \mathbb{Q}(\zeta_m, \zeta_{mn}) \) and \( \mathbb{Q}(\zeta_m) \cap \mathbb{Q}(\zeta_n) = \mathbb{Q} \).

Proof. As \( m \) and \( n \) are relatively prime, \( \zeta_m \zeta_n \) is a primitive \( (mn) \text{th} \) root of unity. Since the Euler phi-function is multiplicative, the result follows from the theorem. \( \Box \)

By the results in §E, we know \( G(\mathbb{Q}(\zeta_n)/\mathbb{Q}) \). Indeed if \( n = 2^e p_1^{e_1} \cdots p_r^{e_r} \) is a factorization with \( 2 < p_1 < \cdots < p_r \) primes, with \( e \geq 0 \) and \( e_i > 0 \) for all \( i \), then

\[
G(\mathbb{Q}(\zeta_n)/\mathbb{Q}) \cong (\mathbb{Z}/2^e\mathbb{Z})^\times \times (\mathbb{Z}/p_1^{e_1}\mathbb{Z})^\times \times \cdots \times (\mathbb{Z}/p_r^{e_r}\mathbb{Z})^\times
\]

by the Chinese Remainder Theorem. We know that \( (\mathbb{Z}/p_i^{e_i}\mathbb{Z})^\times \) is cyclic of order \( \varphi(p_i^{e_i}) = p_i^{e_i-1}(p_i - 1) \) for every \( i \) and if \( n \) is even \( (\mathbb{Z}/2^e\mathbb{Z})^\times \) is cyclic only if \( e = 1, 2 \), otherwise \( (\mathbb{Z}/2^e\mathbb{Z})^\times \cong (\mathbb{Z}/2\mathbb{Z}) \times (\mathbb{Z}/2^{e-2}\mathbb{Z}) \) if \( e \geq 3 \).

We next wish to prove a special case of the following well-known result that uses the same idea as the proof of the above as well as our previous idea in proving the infinitude of primes.

Theorem 58.9. (Dirichlet’s Theorem on Primes in an Arithmetic Progression) Let \( m \) and \( n \) be relatively prime integers. Then there exist infinitely many primes \( p \) satisfying \( p \equiv m \mod n \).

We shall not prove this general theorem. The standard proof uses complex analysis, although there is an elementary proof, i.e., one that uses no complex analysis. The special case that we shall prove is the following:

Proposition 58.10. Let \( n > 1 \) be an integer. Then there exist infinitely many primes \( p \) satisfying \( p \equiv 1 \mod n \).

To prove this, we first establish the following lemma:

Lemma 58.11. Let \( a \) be a positive integer and \( p \) a (positive) prime not dividing \( n \). If \( p \mid \Phi_n(a) \) in \( \mathbb{Z} \), then \( p \equiv 1 \mod n \).

Proof. Let \( - : \mathbb{Z} \rightarrow \mathbb{Z}/p\mathbb{Z} \) be the canonical epimorphism. We know that \( t^n - 1 = \prod_{d \mid n} \Phi_d \) is a factorization of \( t^n - 1 \) into irreducible polynomials in \( \mathbb{Z}[t] \). In particular, \( a^n \equiv 1 \mod p \). Let \( m \) be the order of \( a \) in \( \mathbb{Z}/p\mathbb{Z} \). So \( m \mid n \). If \( m < n \), then

\[
t^n - 1 = \prod_{d \mid n} \Phi_d = \Phi_n \prod_{d \mid n} \Phi_d = \Phi_n(t^m - 1) \prod_{d \mid m} \Phi_d.
\]
It follows that $\bar{a}$ is a multiple root of $t^n - 1$ in $\mathbb{Z}/p\mathbb{Z}$. But $p \not| n$, so $t^n - 1$ has no multiple roots in $\mathbb{Z}/p\mathbb{Z}$. It follows that $m = n$, i.e., $n$ is the order of $\bar{a}$, so $n | p - 1$ by Fermat’s Little Theorem. The result follows.

□

Proof. (of the Proposition) Suppose that the result is false, and $p_1, \ldots, p_r$ are all the primes $p$ satisfying $p \equiv 1 \mod n$. As $\Phi_n$ is monic, we can choose an integer $N > 1$ satisfying $\Phi_n(M) > 1$ for all integers $M > N$ (by calculus). Let $M = p_1 \cdots p_r n N > N$ and $p$ a prime such that $p | \Phi_n(p_1 \cdots p_r n N) > 1$. (Replace $p_1 \cdots p_r$ by 1 if $r = 0$.) Since $t^n - 1 = \prod_{d|n} \Phi_d$ in $\mathbb{Z}[t]$, the constant term of $\Phi_d$ is $\pm 1$, so

$$\Phi_n(p_1 \cdots p_r n N) \equiv \begin{cases} \pm 1 \mod p_i & \text{for } 1 \leq i \leq r \\ \pm 1 \mod n. \end{cases}$$

In particular, $p \neq p_1, \ldots, p_r$ and $p \not| n$. It follows that $p \equiv 1 \mod n$ by the lemma, a contradiction.

□

The proof of the full theorem shows that primes are “equally distributed” among the integers $m$ relatively prime to $n$ with $1 \leq m \leq n$.

We next turn to a special case of another important theorem whose proof we shall also omit, viz.,

**Theorem 58.12.** (Kronecker-Weber Theorem) Let $K/\mathbb{Q}$ be an abelian extension of fields. Then there exists a root of unity $\zeta$ in $\mathbb{C}$ such that $K$ is a subfield of $\mathbb{Q}(\zeta)$

We shall prove this in the case that $K$ is a quadratic extension of $\mathbb{Q}$. The general theorem shows that abelian extensions of $\mathbb{Q}$ are determined by the unit circle. A similar geometric interpretation is true when $\mathbb{Q}$ is replaced by an imaginary quadratic extension of $\mathbb{Q}$, i.e., by $\mathbb{Q}(\sqrt{-d})$, with $d = 1$ or $d > 1$ a square-free integer, and the circle replaced by an appropriate “elliptic curve”. Unfortunately, this does not extend further, although a major triumph in number theory was the determination of abelian extensions $L$ of $K$ when $K/\mathbb{Q}$ is finite. There just is no geometric formulation in this general case.

If $a$ is an integer not divisible by $p$, recall we defined the Legendre symbol in Definition 31.8 to be

$$\left( \frac{a}{p} \right) := \begin{cases} +1 & \text{if } a \mod p \text{ is a square}. \\ -1 & \text{if } a \mod p \text{ is not a square}. \end{cases}$$

It is convenient to define $\left( \frac{a}{p} \right) = 0$ if $p \not| a$. Clearly, the Legendre symbol $\left( \frac{a}{p} \right)$ only depends on $a \mod p$. Let $p$ be any prime and consider the squaring map $f : (\mathbb{Z}/p\mathbb{Z})^\times \rightarrow (\mathbb{Z}/p\mathbb{Z})^\times$ given by $x \mapsto x^2$. Then $f$
is a group homomorphism with \( \ker f = \{ \pm 1 \} \) with image \( (\mathbb{Z}/p\mathbb{Z})^\times \).
In particular, if \( p \) is odd, then \( |\text{im } f| = (p - 1)/2 \), i.e., half of \( (\mathbb{Z}/p\mathbb{Z})^\times \) are squares and half non-squares.

**Lemma 58.13.** Let \( p \) be an odd prime. Then for all integers \( a \) and \( b \) not divisible by \( p \), we have the following:

1. \( \left( \frac{a}{p} \right) = \left( \frac{b}{p} \right) \) if \( a \equiv b \mod p \).
2. \( \left( \frac{ab}{p} \right) = \left( \frac{a}{p} \right) \left( \frac{b}{p} \right) \).
3. (Euler’s Criterion) If \( p \nmid a \), then \( \left( \frac{a}{p} \right) \equiv (a)^{(p-1)/2} \mod p \). In particular, \( \left( \frac{-1}{p} \right) = (-1)^{(p-1)/2} \).
4. \( \sum_{a=1}^{p-1} \left( \frac{a}{p} \right) = 0 \).

**Proof.** We have already observed (1), and (4) is the fact that half of \( (\mathbb{Z}/p\mathbb{Z})^\times \) are squares and half not.

(3): We know that \( x^{p-1} \equiv 1 \mod p \) if \( p \nmid x \), so if \( a \) is a square modulo \( p \), then \( a^{(p-1)/2} \equiv 1 \mod p \). Conversely, suppose that \( a \) is a square modulo \( p \). As \( (\mathbb{Z}/p\mathbb{Z})^\times \) is a cyclic group of order \( p - 1 \), we have \( (\mathbb{Z}/p\mathbb{Z})^\times = \langle x \rangle \) for some \( x \). Suppose that \( a \equiv x^n \mod p \). Then \( n \) must be even lest \( x^{(p-1)/2} \equiv 1 \mod p \), hence \( a \) is a square modulo \( p \). The case of \( a = -1 \) follows easily as \( \left( \frac{-1}{p} \right) - (-1)^{(p-1)/2} \) is \(-2, 0, \) or \(2 \). (In fact, we have previously shown this case.)

(2) follows from (3). \( \square \)

**Proposition 58.14.** Let \( p \) be an odd prime and \( \zeta \) a primitive \( p \)th root of unity. Set

\[
S := \sum_{a=1}^{p-1} \left( \frac{a}{p} \right) \zeta^a
\]

in \( \mathbb{Q}(\zeta) \) (even in \( \mathbb{Z}[\zeta] \)). Then

\[
S^2 = \left( \frac{-1}{p} \right) p = (-1)^{p-1}/2 \mod p \text{ in } \mathbb{Q}(\zeta) \text{ (even in } \mathbb{Z} \).
\]

In particular,

- \( \sqrt{p} \) lies in \( \mathbb{Q}(\zeta) \) if \( p \equiv 1 \mod 4 \)
- \( \sqrt{-p} \) lies in \( \mathbb{Q}(\zeta) \) if \( p \equiv 3 \mod 4 \).

Moreover, \( \mathbb{Q}(S) \) is the unique quadratic extension of \( \mathbb{Q} \) lying in \( \mathbb{Q}(\zeta) \).
The sum $S$ in the proposition is called a Gauss sum.

**Proof.** We have

$$S^2 = \sum_{a=1}^{p-1} \sum_{b=1}^{p-1} \left( \frac{a}{p} \right) \left( \frac{b}{p} \right) \zeta^{a+b} = \sum_{a=1}^{p-1} \left( \frac{ab}{p} \right) \zeta^{a+b}.$$

As $a$ ranges over $1, \ldots, p-1$ mod $p$, so does $ab$ for $b = 1, \ldots, p-1$ mod $p$, so upon replacing $a$ by $ab$, we see that

$$S^2 = \sum_{a=1}^{p-1} \sum_{b=1}^{p-1} \left( \frac{ab^2}{p} \right) \zeta^{b(a+1)} = \sum_{a=1}^{p-1} \left( \frac{a}{p} \right) \zeta^{b(a+1)}$$

(*)

where the first term on the second line arises from the term $a = p-1$. As $\zeta^r$ is also a primitive $p$th root of unity for $r = 1, \ldots, p-1$, we have

$$1 + \zeta^r + (\zeta^r)^2 + \cdots + (\zeta^r)^{p-1} = 0$$

for such $r$. In particular, we have

$$\sum_{b=1}^{p-1} \zeta^{b(a+1)} = -1$$

if $a \neq p-1$ in (*). Therefore, we have, using Lemma 58.13(4),

$$S^2 = \sum_{b=1}^{p-1} \left( \frac{-1}{p} \right) - \sum_{a=1}^{p-2} \left( \frac{a}{p} \right) = (p-1)\left( \frac{-1}{p} \right) + \left( \frac{p-1}{p} \right) = p\left( \frac{-1}{p} \right).$$

Consequently, as $S \subset \mathbb{Q}(\zeta)$, we have $\sqrt{p\left( \frac{-1}{p} \right)}$ lies in $\mathbb{Q}(\zeta)$. The “in particular statement” follows from Euler’s Criterion. Finally as $G(\mathbb{Q}(\zeta)/\mathbb{Q}) \cong (\mathbb{Z}/p\mathbb{Z})^\times$ is a cyclic group of order $p-1$, there exists a unique subgroup $H$ of $G(\mathbb{Q}(\zeta)/\mathbb{Q})$ of index two. By the Fundamental Theorem of Galois Theory, we must have $\mathbb{Q}(S) = \mathbb{Q}(\zeta)^H$.

**Remark 58.15.** Let $p$ be an odd prime and $\zeta$ a primitive $p^r$th root of unity in $\mathbb{C}$. Then $\mathbb{Q}(\sqrt{(-1)^{p-1}p})$ is the unique quadratic extension of $\mathbb{Q}$ in $\mathbb{Q}(\zeta)$ by Remark 58.7.

**Corollary 58.16.** Let $p$ be a (positive) prime, $\zeta_p$ a primitive $p$th root of unity, and $\zeta_{4p}$ a primitive $4p$th root of unity. Then $\sqrt{p}$ and $\sqrt{-p}$ lie in $\mathbb{Q}(\zeta_{4p})$.

**Proof.** Let $\zeta_n$ be a fixed primitive $n$th root of unity in $\mathbb{C}$ for each $n$.

**Case 1.** $p$ an odd prime:

If $p \equiv 1 \mod 4$, then by the proposition, we know that $\sqrt{p}$ lies in $\mathbb{Q}(\zeta_p) \subset \mathbb{Q}(\zeta_{4p})$ and $\sqrt{-p}$ lies in $\mathbb{Q}(\sqrt{-1}, \zeta_p) = \mathbb{Q}(\zeta_{4p})$ by Corollary 58.8. The case that $p \equiv 3 \mod 4$ is analogous.
Case 2. $p = 2$:
Let $y = \zeta_8 + \zeta_8^{-1}$. Then $\zeta_8$ is the point $(1 + \sqrt{-1})/2$ on the unit circle in the complex plane, i.e., the intersection with the $45^\circ$ ray with the real axis and $\zeta_8^{-1}$ is its complex conjugate. Computation shows $y^2 = 2$. It follows that $\sqrt{2}$ lies in $\mathbb{Q}(\zeta_8)$. As $\sqrt{-1}$ lies in $\mathbb{Q}(\zeta_4) \subset \mathbb{Q}(\zeta_8)$, it follows that $\sqrt{2}$ and $\sqrt{-2}$ lie in $\mathbb{Q}(\zeta_8)$. □

Theorem 58.17. Let $n$ be a nonzero integer. Then $\sqrt{n}$ lies in $\mathbb{Q}(\zeta_{4n})$, where $\zeta_{4n}$ is a primitive $4n$th root of unity. In particular, any quadratic extension of $\mathbb{Q}$ lies in $\mathbb{Q}(\zeta_r)$ for some root of unity $\zeta_r$ in $\mathbb{C}$.

Proof. Let $\zeta_n$ be a fixed primitive $n$th root of unity in $\mathbb{C}$ for each $n$. We may assume that $n$ is square-free. We know the result if $n = -1$, so we may assume that $n = \pm p_1 \cdots p_r$ with $p_1 < \cdots < p_r$ positive primes with $r \geq 1$. Then $\sqrt{n}$ is an element of $\mathbb{Q}(\zeta_4, \zeta_{p_1}, \ldots, \zeta_{p_r}) = \mathbb{Q}(\zeta_{4n})$ if $n$ is odd and $\sqrt{n}$ is an element of $\mathbb{Q}(\zeta_8, \zeta_{p_2}, \ldots, \zeta_{p_r}) = \mathbb{Q}(\zeta_{4n})$ if $n$ is even. □

We turn to a proof of Quadratic Reciprocity, one of hundreds!

Theorem 58.18. (Law of Quadratic Reciprocity) Let $p$ and $q$ be distinct odd primes. Then

\[
\left( \frac{p}{q} \right) \left( \frac{q}{p} \right) = (-1)^{\frac{p-1}{2} \frac{q-1}{2}}.
\]

In addition, we have

(Supplement #1) \( \left( \frac{-1}{p} \right) = (-1)^{\frac{p-1}{2}}. \)

(Supplement #2) \( \left( \frac{2}{p} \right) = (-1)^{\frac{p^2-1}{8}}. \)

Proof. Let $\zeta$ be a primitive $p$th root of unity in $\mathbb{C}$ and

\[
S := \sum_{a=1}^{p-1} \left( \frac{a}{p} \right) \zeta^a \text{ in } \mathbb{Z}[\zeta].
\]

By Proposition 58.14 and Euler’s Criterion, we have

\[
S^q = S(S^2)^{\frac{q-1}{2}} = S \left( \frac{1}{p} \right)^{\frac{q-1}{2}} p^{\frac{q-1}{2}} = S(-1)^{\frac{p-1}{2} \frac{q-1}{2}} p^{\frac{q-1}{2}}.
\]

Multiplying this equation by $S$, we see that

\[
S^{q+1} = S^2(S^2)^{\frac{q-1}{2}} = S^2(-1)^{\frac{p-1}{2} \frac{q-1}{2}} p^{\frac{q-1}{2}}
\]

\[=(*) \]

\[
p^{\frac{q-1}{2}} \left( \frac{-1}{p} \right) (-1)^{\frac{p-1}{2} \frac{q-1}{2}} p^{\frac{q-1}{2}}.
\]
hence is an integer. By the Binomial Theorem (actually the Multinomial Theorem), since $q$ is an odd prime, we see upon multiplying out and coalescing terms that

$$S^q = \left( \sum_{a=1}^{p-1} \left( \frac{a}{p} \right) \zeta^a \right)^q = \sum_{a=1}^{p-1} \left( \frac{a}{p} \right) q \zeta^a + qr = \sum_{a=1}^{p-1} \left( \frac{a}{p} \right) \zeta^a + qr$$

for some $r$ in $\mathbb{Z}[\zeta]$. As $aq, a = 1, \ldots, p - 1$, runs over all the nonzero residue classes modulo $p$, we have

$$S^q = \sum_{a=1}^{p-1} \left( \frac{q^2 a}{p} \right) \zeta^{aq} + qr = \sum_{a=1}^{p-1} \left( \frac{aq}{p} \right) \zeta^{aq} + qr$$

Multiplying this by $S$ then yields

$$S^{q+1} = \left( \frac{q}{p} \right) S^2 + qr = \left( \frac{q}{p} \right) (-1)^{p-1} p + qr S$$

in $\mathbb{Z}[\zeta]$.

**Check 58.19.** $q\mathbb{Z}[\zeta_p] \cap \mathbb{Z} = q\mathbb{Z}$.

So we have, using (*),

$$qrS = S^{q+1} - \left( \frac{q}{p} \right) (-1)^{p-1} p \text{ lies in } q\mathbb{Z},$$

which implies

$$S^{q+1} \equiv \left( \frac{q}{p} \right) (-1)^{p-1} p \mod q.$$ 

Using (*) again, now shows that

$$p \left( \frac{-1}{p} \right) (-1)^{p-1} \frac{q}{p} \frac{q-1}{p} \equiv \left( \frac{q}{p} \right) (-1)^{p-1} p \mod q.$$ 

Since $p \mod q$ is a unit in $\mathbb{Z}/q\mathbb{Z}$ and $p^{\frac{q-1}{2}} \equiv \left( \frac{p}{q} \right) \mod q$ by Euler’s Criterion, we have

$$\left( \frac{p}{q} \right) = (-1)^{\frac{p-1}{2} \frac{q-1}{2}} \left( \frac{q}{p} \right) \mod q.$$ 

Finally, as $\left( \frac{p}{q} \right) - (-1)^{\frac{p-1}{2} \frac{q-1}{2}} \left( \frac{q}{p} \right)$ is $-2, 0, \text{ or } 2$, the result follows.

Since we have already proven Supplements #1, to finish we need only establish Supplement #2. We work in $\mathbb{Z}/p\mathbb{Z}$ with $p$ an odd prime. Let $\zeta$ be a primitive 8th root of unity in an algebraic extension of $\mathbb{Z}/p\mathbb{Z}$, i.e., a root of $t^8 - 1$ over $\mathbb{Z}/p\mathbb{Z}$ generating the 8th roots of unity and
set $y = \zeta + \zeta^{-1}$ in $(\mathbb{Z}/p\mathbb{Z})(\zeta)$. As $\zeta^8 = 1$ and $\zeta$ has order 8, we have $\zeta^4 = -1$ in $(\mathbb{Z}/p\mathbb{Z})(\zeta)$ with $1 \neq -1$ since $p$ is odd. Multiplying this last equation by $\zeta^{-2}$, we see that $\zeta^2 + \zeta^{-2} = 0$ in $(\mathbb{Z}/p\mathbb{Z})(\zeta)$, hence $y^2 = (\zeta + \zeta^{-1})^2 = \zeta^2 + 2 + \zeta^{-2} = 2$. [Cf. this to the case in characteristic zero.] Let $\zeta: \mathbb{Z} \to \mathbb{Z}/p\mathbb{Z}$ be the canonical epimorphism. By Euler’s Criterion, $(\frac{2}{p}) = 2^{p-1} \equiv y^p - y \pmod{p}$. [Note that $y \notin \mathbb{Z}/p\mathbb{Z}$, so $y^p - y$ is not necessarily one.] By the Children’s Binomial Theorem, $y^p = \zeta^p + \zeta^{-p}$ in $(\mathbb{Z}/p\mathbb{Z})(\zeta)$.

**Case 1.** $p \equiv \pm 1 \pmod{8}$:

As $\zeta^8 = 1$, we have

$$y^p = \zeta^p + \zeta^{-p} = \zeta + \zeta^{-1} = y \pmod{\mathbb{Z}/p\mathbb{Z}}(\zeta),$$

so $(\frac{2}{p}) = y^{p-1} = 1$ in $(\mathbb{Z}/p\mathbb{Z})(\zeta)$, hence $(\frac{2}{p}) = 1$ in $\mathbb{Z}$.

**Case 2.** $p \equiv \pm 5 \pmod{8}$:

As $\zeta^4 = -1$, we have

$$y^p = \zeta^p + \zeta^{-p} = \zeta^5 + \zeta^{-5} = -\zeta - \zeta^{-1} = -y \pmod{\mathbb{Z}/p\mathbb{Z}}(\zeta),$$

so $(\frac{2}{p}) = y^{p-1} = -1$ in $(\mathbb{Z}/p\mathbb{Z})(\zeta)$, hence $(\frac{2}{p}) = -1$ in $\mathbb{Z}$.

This completes the proof. □

**Examples 58.20.**

1. Quadratic Reciprocity and the properties of the Legendre symbol allows computing Legendre symbols quite efficiently. For example,

$$\left(\frac{29}{43}\right) = (-1)^{29-1} \frac{43-1}{2} \left(\frac{43}{29}\right) = \left(\frac{43}{29}\right) = \left(\frac{14}{29}\right) = \left(\frac{2}{29}\right) \left(\frac{7}{29}\right)$$

$$= (-1)^{\frac{(29-1)(29+1)}{8}} \left(\frac{7}{29}\right) = -\left(\frac{7}{29}\right) = (-1)^{\frac{29-1}{2}} \left(\frac{29}{7}\right)$$

$$= -\left(\frac{29}{7}\right) = -\left(\frac{1}{7}\right).$$

So 29 is not a square modulo 43.

2. The Legendre Symbol can also be use to study Diophantine equations. Let $x$ and $y$ be variables. We show that the Diophantine equation $y^2 = x^3 + 45$ has no solution, i.e., no solution in integers. Suppose that $(x_0, y_0) \in \mathbb{Z}^2$ is a solution. We first show that we cannot have $3 \mid x_0$. If $x_0 = 3x_1$, then $y_0 = 3y_1$ for some integers $x_1, y_1$. It follows that $3^2y_1^2 = 3^3x_1^3 + 3^25$ or $y_1^2 \equiv 5 \equiv 2 \pmod{3}$ which is impossible. If $x_0$ is even, then $y_0^2 \equiv 45 \equiv 5 \pmod{8}$ which is impossible and if $x_0 \equiv 1 \pmod{4}$, then $y_0^2 \equiv 46 \equiv 2 \pmod{4}$, which is also
impossible. Therefore, we are reduced to the case \( x_0 \equiv 3 \mod 4 \), equivalently, \( x_0 \equiv 3, 7 \mod 8 \).

**Case.** \( x_0 \equiv 7 \equiv -1 \mod 8 \).

As \( y_0^2 = x_0^3 + 45 \), we have \( y_0^2 - 2 \cdot 3^2 \equiv x_0^3 + 27 = (x_0 + 3)(x_0^2 - 3x + 9) \) and \( x_0^2 - 3x_0 + 9 \equiv -3 \mod 8 \). Therefore, there exists a prime \( p \) satisfying \( p \equiv \pm 3 \mod 8 \) and \( p \mid x_0^2 - 3x_0 + 9 \). This implies that \( y_0^2 = 2 \cdot 3^2 \mod p \). If \( p \mid y_0 \), then \( p \mid x_0 \) which is impossible. So \( p \not\equiv 3 \), hence \( \left( \frac{2 \cdot 3^2}{p} \right) = \left( \frac{2}{p} \right) = -1 \), contradicting the Second Supplement.

**Case.** \( x_0 \equiv 3 \mod 8 \).

We have \( y_0^2 - 2 \cdot 6^2 = x_0^3 - 27 = (x_0 - 3)(x_0^2 + 3x_0 + 9) \), so there exists a prime \( p \), satisfying \( p \equiv \pm 3 \mod 8 \) and \( p \mid x_0^2 + 3x_0 + 9 \), as \( x_0^2 + 3x_0 + 9 \equiv 3 \mod 8 \). Therefore, \( y_0^2 \equiv 2 \cdot 6^2 \mod p \) has a solution. Since \( p \not\equiv 3 \), hence \( \left( \frac{2 \cdot 6^2}{p} \right) = \left( \frac{2}{p} \right) = -1 \), contradicting the Second Supplement.

Let \( K/\mathbb{Q} \) be a finite field extension. Number theory studies the set \( \mathcal{O}_K := \{ x \in K \mid x \text{ is a root of a monic polynomial } f \in \mathbb{Z}[x] \} \). We mention some of the facts about this here, and study this subject in Chapter XV. The set \( \mathcal{O}_K \) is a ring, hence a domain. It is called the ring of algebraic integers in \( K \). Although it is not a UFD, it has an analogous property about ideals, viz, that every proper ideal in \( \mathcal{O}_K \) is a product of prime ideals, unique up to order. Such a domain is called a Dedekind domain. It turns out that every nonzero prime ideal \( \mathfrak{p} \) in \( \mathcal{O}_K \) is maximal, so \( \mathcal{O}_K / \mathfrak{p} \) is a field, in fact a finite field so a finite extension of \( \mathbb{Z}/p\mathbb{Z} \) where \( \mathfrak{p} \cap \mathbb{Z} = p\mathbb{Z} \) for some prime \( p \). Let \( f_\mathfrak{p} \) denote the degree of this extension. One studies the factorization of \( p\mathbb{Z} \), with \( p \) a prime integer, in \( \mathcal{O}_K \). Suppose that this factorization is \( p\mathcal{O}_K = \mathfrak{p}_1^{e_1} \cdots \mathfrak{p}_r^{e_r} \), with \( \mathfrak{p}_i \) distinct prime ideals in \( \mathcal{O}_K \) and \( e_i \) positive integers. A basic result is that \( [K : \mathbb{Q}] = \sum e_i f_\mathfrak{p}_i \). If, in addition, \( K/\mathbb{Q} \) is Galois then all the \( e_i \)'s are equal, say equal \( e \), and all the \( \mathfrak{p}_i \)'s are equal, say equal \( f \), so \( [K : \mathbb{Q}] = e f r \). The Legendre symbol provides information about the case that \( K/\mathbb{Q} \) is a quadratic extension (hence Galois), say \( K = \mathbb{Q}(\sqrt{d}) \) with \( d \) a square-free integer or \( -1 \). In this case \( \mathcal{O}_K = \mathbb{Z}[D] \) where \( D = d \) if \( d \equiv 2 \) or \( 3 \mod 4 \) and \( D = (-1 + \sqrt{d})/2 \) if \( d \equiv 1 \mod 4 \). Let \( p \) be a prime integer. Three possibilities can occur (called the splitting type of \( p \)).

1. \( r = 2 \), i.e., \( p\mathcal{O}_K \) factors as a product of two distinct primes – we say that \( p \) splits completely in \( \mathcal{O}_K \).
(ii) $f = 1$, i.e., $p\mathcal{O}_K$ is a prime ideal in $\mathcal{O}_K$ – we say that $p$ is inert in $\mathcal{O}_K$.

(iii) $e = 1$, i.e., $p\mathcal{O}_K$ is the square of a prime ideal in $\mathcal{O}_K$ – we say that $p$ ramifies in $\mathcal{O}_K$.

A prime $p$ ramifies in $\mathcal{O}_K$ if and only if $p \mid d_K$, where $d_K = 4D$ if $D \equiv 2$ or $3 \pmod{4}$ and $d_K = D$ if $d \equiv 1 \pmod{4}$. Let $p$ be an odd prime. Then $p$ ramifies if and only if $p \mid D$, so assume this is not the case. Then $f = 1$ or $2$ depending on whether $d$ is a square or not modulo $p$, i.e., on $\left(\frac{d}{p}\right)$. If $p\mathbb{Z}$ splits completely and $a$ is a square modulo $p$, then $p\mathcal{O}_K = (p, a + \sqrt{d})(p, a - \sqrt{d})$. Note that if $p$ and $q$ are odd primes, the Law of Quadratic Reciprocity tells the relation of the splitting type of $p$ in $\mathcal{O}_{\mathbb{Z}[\sqrt{q}]}$ and the splitting type of $q$ in $\mathcal{O}_{\mathbb{Z}[\sqrt{p}]}$. The splitting type of $2$ is more complicated.

**Exercises 58.21.**

1. Show that the set of all arithmetic functions satisfies the associative law under the Dirichlet product. In particular, the set $\mathcal{A}$ of arithmetic functions nonzero at 1 is an abelian monoid.

2. Let $N$, $\tau$, and $\sigma$ be the arithmetic function functions defined by

$$N(n) = n,$$

$$\tau(n) = \text{the number of divisors of } n, \text{ i.e., } \tau(n) = \sum_{d \mid n} 1,$$

$$\sigma(n) = \text{the sum of the divisors of } n, \text{ i.e., } \sigma(n) = \sum_{d \mid n} d$$

for all $n$, respectively. Show that $N = \varphi \ast U$, $\varphi = N \ast \mu$, $\sigma = \varphi \ast \tau$, and $\varphi \ast \sigma = N \ast N$.

3. Let $\zeta(s) := \sum \frac{1}{n^s}$, $s$ a real (or complex) variable. Show the following:

$$\frac{\zeta(s - 1)}{\zeta(s)} = \sum_{1}^{\infty} \frac{\varphi(n)}{n^s} \text{ if (the real part of) } s > 2.$$  

$$\zeta(s)^2 = \sum_{1}^{\infty} \frac{\tau(n)}{n^s} \text{ if (the real part of) } s > 1.$$  

$$\zeta(s - 1)\zeta(s) = \sum_{1}^{\infty} \frac{\sigma(n)}{n^s} \text{ if (the real part of) } s > 2.$$  

4. Show all of the following:

(i) The Dirichlet product of two multiplicative functions is multiplicative.
(ii) If $f$ and $g$ are arithmetic functions not zero at one with $g$ and $f \ast g$ multiplicative, then $f$ is multiplicative.

5. The set of multiplicative functions is an abelian group under the Dirichlet product.

6. Prove the following versions of the M"obius Inversion Theorem.
   (a) If $f, g : \mathbb{Z}^+ \to G$ with $G$ an additive group, then
   
   $$f(n) = \sum_{d \mid n} g(d) \text{ if and only if } g(n) = \sum_{d \mid n} f(d)\mu\left(\frac{n}{d}\right)$$
   
   for all $n \in \mathbb{Z}^+$, where $\mu$ is viewed as a function to $G$.
   (b) If $f, g : \mathbb{Z}^+ \to G$ with $G$ a multiplicative group, then
   
   $$f(n) = \prod_{d \mid n} g(d) \text{ if and only if } g(n) = \prod_{d \mid n} f(d)\mu\left(\frac{n}{d}\right)$$
   
   for all $n \in \mathbb{Z}^+$, where $\mu$ is viewed as a function to $G$.

7. Let $F$ be any field and $d$ a positive integer. Show that $t^{d} - 1 \mid t^{n} - 1$ in $F[t]$ if and only if $d \mid n$.

8. Show that every finite cyclic group occurs as a Galois group $G(K/Q)$ for some Galois extension $K/Q$.

9. Prove that every finite abelian group occurs as a Galois group over the rational numbers. (Hint: Use Proposition 58.10.)


11. Let $K = \mathbb{Q}(\sqrt{d})$ with $d$ a square-free integer (or $-1$) and $p$ an odd prime integer. Show directly that $p \mathcal{O}_K = (p, a + \sqrt{d})(p, a - \sqrt{d})$ in $\mathcal{O}_K$ if $p$ does not divide $d$ and $a$ in $\mathbb{Z}$ is a square modulo $p$.

12. Let $K = \mathbb{Q}(\sqrt{d})$ with $d$ a square-free integer (or $-1$) and $p$ an odd prime integer. Show directly that $p \mathcal{O}_K$ is a prime ideal if $\mathcal{O}_K$ if $a$ in $\mathbb{Z}$ is not a square modulo $p$.

13. Let $K = \mathbb{Q}(\sqrt{d})$ with $d$ a square-free integer (or $-1$) and $p$ an odd prime integer. Show directly that $p \mathcal{O}_K = (p, \sqrt{d})^2$ in $\mathcal{O}_K$ if $p \mid d$.

59. Radical Extensions

In this section, we generalize the notion of square root towers and solve the problem of when a formula for roots of a polynomial involving only addition, multiplication (in which we include taking inverses) and extraction of $n$th roots of elements for various $n$ over the rational numbers exists. The key is if the Galois group of a polynomial is a solvable group.
Definition 59.1. An extension of fields $K/F$ is called a radical extension if there exist elements $u_i$ in $K$ for $1 \leq i \leq m$ (some $m$) such that $K = F(u_1, \ldots, u_m)$ and for each $i$ there exists a positive integer $n_i$ such that $u_i^{n_i}$ lies in $F$ and $u_i^{n_i} \in F(u_1, \ldots, u_{i-1})$ for each $i > 1$.

As one would expect, radical extensions have properties similar to square root towers, and they do.

Remark 59.2. Let $K/F$ and $L/F$ be field extensions with $K$ and $L$ lying in an extension field of $K$.
1. If $K/F$ is a square root tower, then $K/F$ is radical.
2. If $K/F$ is the splitting field of $t^n - 1$ in $F[t]$, then $K/F$ is radical.
3. If $K/F$ is the splitting field of $t^n - a$ in $F[t]$, then $K/F$ is radical, as $K = F(\omega, \theta)$ with $\omega$ a primitive $n$th root of unity in $K$ and $\theta$ an element in $K$ such that $\theta^n = a$ in $K$.
4. If $K/F$ is radical, then so is $L(K)/L$.
5. If $L/K$ and $K/F$ are both radical, then so is $L/F$.
6. If $L/F$ is radical and $L/K$, then $L/K$ is radical.
7. If $E/L_i/F$ are intermediate fields with each $L_i/F$ radical for $i = 1, \ldots, n$, then $F(L_1 \cup \cdots \cup L_n)/F$ is radical.
8. If $K/F$ is radical and $\sigma : K \to L$ is a (field) homomorphism, then $\sigma(K)/\sigma(F)$ is radical.
9. If $K/F$ is radical and $L/K$ is a normal closure of $K/F$, then $L/F$ is radical. Indeed this follows from the last two remarks, as $L = F(\cup_{G(L/F)} \sigma(K))$.
10. If $K/F$ is a square root tower and $L/K$ is a normal closure of $K/F$, then $L/F$ is a square root tower.

The last remark allows us to refine the Constructibility Criterion 51.9.

Theorem 59.3. (Constructibility Criterion (Refined Form)) Let $z$ be a complex number and $z_1(= 0), z_2(= 1), \ldots, z_n$ other complex numbers with $n \geq 2$. Set $F = \mathbb{Q}(z_1, z_2, \ldots, z_n, \bar{z}_1, \ldots, \bar{z}_n)$. Then the following are equivalent:

1. The complex number $z$ is constructible from $z_1, \ldots, z_n$.
2. The complex number $z$ is algebraic over $F$ and the normal closure of $F(z)/F$ in $\mathbb{C}$ is a square root tower.
3. The complex number $z$ is algebraic over $F$ and if $E/F(z)$ is the normal closure of $F(z)/F$ in $\mathbb{C}$, then $[E : F] = 2^e$ for some $e$. 


Proof. (2) if and only if (3) follows from the Square Root Tower Theorem 56.11.
(2) ⇒ (1) follows from the original Constructibility Criterion 51.9.
(1) ⇒ (2): By the original Constructibility Criterion 51.9, there exists a square root tower $K/F$ with $z \in K$. By Remark (10) above, we may assume that $K/F$ is normal. Let $K/E/F(z)/F$ with $E/F(z)$ the normal closure of $F(z)/F$ in $K$. Since $[E : F] | [K : F] = 2^e$, some $e$, we are done by the Square Root Tower Theorem 56.11.

We wish to find the group theoretic criterion for a finite extension to be a radical extension. For simplicity we shall establish this in the case that the ground field is of characteristic zero.

Let $K/F$ be a finite Galois extension of fields. Recall that we call it an abelian extension (respectively, cyclic extension) if $G(K/F)$ is abelian (respectively, cyclic).

Example 59.4. 1. Any cyclic extension of fields is abelian.
2. Suppose $n$ is a positive integer and either $\text{char } F = 0$ or $\text{char } F \nmid n$. If $K$ is a splitting field of $t^n - 1$ over $F$, then $K/F$ is abelian (and, in fact, cyclic if $n$ is $2, 4, p^r, 2p^r$, with $p$ an odd prime, cf. E.6 below).
3. If $K/F$ is abelian (respectively, cyclic) and $K/E/F$ is an intermediate field, then both $K/E$ and $E/F$ are abelian (respectively, cyclic).
4. Suppose $n$ is a positive integer and either $\text{char } F = 0$ or $\text{char } F \nmid n$. Suppose that $t^n - 1$ splits over $F$. If $K$ is a splitting field of $t^n - a$ over $F$, then $K/F$ is cyclic.

We review our study about solvable groups.

Review 59.5. A group $G$ is called solvable (respectively, polycyclic) if there exist a finite sequence of subgroups of $G$:

$$1 = N_0 \subset N_1 \subset \cdots \subset N_r = G$$

satisfying $N_i \triangleleft N_{i+1}$ and $N_{i+1}/N_i$ is abelian (respectively, cyclic) for every $i = 1, \ldots, r - 1$, i.e., $G$ has a finite subnormal series with abelian (respectively cyclic) factors. Let $G$ be a group. We know the following facts about solvable and polycyclic groups:

1. Abelian groups are solvable.
2. If $G$ is solvable and $H$ a subgroup of $G$, then $H$ is solvable.
3. If $N$ is a normal subgroup of $G$, then $G$ is solvable if and only if both $N$ and $G/N$ are solvable.
4. If $G$ is polycyclic, then $G$ is solvable.
5. If $G$ is finite and solvable, then $G$ is polycyclic.
6. The alternating group $A_n$ with $n \geq 5$ is a nonabelian simple group, hence not solvable.

7. The symmetric group $S_n$ with $n \geq 5$ is not solvable.

Applying the above to field theory we have:

**Remarks 59.6.** Let $K/F$ be a finite extension of fields with $K/E/F$ an intermediate field.

1. If $G(K/F)$ is solvable, then $G(K/E)$ is solvable.

2. If $K/F$ and $E/F$ are Galois, then $G(K/F)$ is solvable if and only if both $G(K/E)$ and $G(E/F)$ are solvable, since, by the Fundamental Theorem of Galois Theory, $G(E/F) \cong G(K/F)/G(K/E)$.

Given a field extension, it is often useful to extend the base field by adjoining appropriate roots of unity. In investigating radical extensions, this is most useful. Indeed using the last remark we easily establish the following lemma.

**Lemma 59.7.** Let $n$ be a positive integer and $F$ a field satisfying $\text{char } F = 0$ or $\text{char } F \nmid n$. Suppose that $K/F$ is a splitting field of a separable polynomial $f$ over $F$ and $L/F$ is a splitting field of $t^n - 1$ over $K$. Then $G(L/F)$ is Galois. Moreover, $G(L/F)$ is solvable if and only if $G(K/F)$ is solvable.

**Proof.** We know that $L$ is a splitting field of the separable polynomial $(t^n - 1)f$ over $F$, so $L/F$ is Galois. As $L/K$ is abelian, by Remark 59.6, the group $G(L/F)$ is solvable if and only if the group $G(K/F)$ is solvable.

The key result is the following theorem.

**Theorem 59.8.** Let $F$ be a field of characteristic zero and $K/F$ a radical extension. If $K/E/F$ is an intermediate field, then $G(E/F)$ is solvable.

**Proof.** We begin with two reductions.

**Reduction 1.** We may assume that $E/F$ is Galois:

Let $F_0 = E^{G(E/F)} \supset F$. Then $K/F_0$ is also radical, $G(E/F) = G(E/F_0)$, and $E/F_0$ is Galois. Replacing $F$ by $F_0$ establishes the reduction.

**Reduction 2.** We may assume that $K/F$ is radical:

Let $L/K$ be a normal closure of $K/F$. Since $K/F$ is radical, so is $L/F$. Replacing $K$ by $L$ establishes the second reduction.
So we may now assume that $K/F$ is a Galois and a radical extension. By Remark 59.6, it now suffices to show that $G(K/F)$ is solvable as $E/F$ is Galois. Therefore, we have reduced to proving the following:

**Claim.** Let $F$ be a field of characteristic zero and $K/F$ a finite extension that is both radical and Galois. Then $G(K/F)$ is solvable:

Let $K = F(u_1, \ldots, u_m)$ and $n_i$ positive integers, $i = 1, \ldots, m$, satisfying $u_1^{n_1}$ lies in $F$ and $u_i^{n_i} \in F(u_1, \ldots, u_{i-1})$ for each $i > 1$. We prove the claim by induction.

**Induction Hypothesis.** Let $L/E$ be a finite, radical, Galois extension of fields with $E$ a field of characteristic zero. Suppose that $L = E(v_1, \ldots, v_{m-1})$ and $r_i$ positive integers, $i = 1, \ldots, m-1$, satisfying $v_1^{r_1}$ lies in $E$ and $v_i^{r_i} \in E(v_1, \ldots, v_{i-1})$ for each $i > 1$. Then $G(L/E)$ is solvable:

Note: The $m = 0$ case is trivial (and the $m = 1$ case is included in the proof below).

With $K = F(u_1, \ldots, u_m)$ as above, let $\zeta$ be a primitive $n_1$st root of unity, so a primitive root of unity of $t^{n_1} - 1$ over $K$. We have the following picture with the explanation to follow:

$$
\begin{array}{c}
K(\zeta) \\
\downarrow \text{gal, rad} \\
K \\
\downarrow \text{gal} \\
F \\
\end{array}
\quad
\begin{array}{c}
F(\zeta) \\
\downarrow \text{cyc} \\
F(\zeta), u_1 \\
\downarrow \text{cyc} \\
F(\zeta). \\
\end{array}
$$

$K(\zeta)/K$ and $F(\zeta)/F$ are Galois and abelian, since a splitting field over $t^{n_1} - 1$ over $K$ and $F$, respectively. The extension $K(\zeta)/F$ is Galois by the lemma, hence $K(\zeta)/F(\zeta)$ and $K(\zeta)/F(\zeta, u_{n_1})$ are Galois. Since $F(\zeta, u_1)$ is a splitting field of $t^{n_1} - u_1^{n_1}$ over $F(\zeta)$, the extension $F(\zeta, u_1)/F$ is Galois, as it is a splitting field of the separable polynomial $(t^{n_1} - u_1^{n_1})(t^{n_1} - 1)$ over $F$. We know that $F(\zeta, u_1)/F(\zeta)$ is cyclic and $F(\zeta)/F$ is abelian, so $G(F(\zeta, u_1)/F)$ is solvable by Remark 59.6.

In particular, all the extensions in the picture above are Galois. We have $K(\zeta) = F(\zeta, u_1)(u_2, \ldots, u_m)$ with $K(\zeta)/F(\zeta, u_1)$ radical and Galois. Consequently, by the induction hypothesis, $G(K(\zeta)/F(\zeta, u_1))$ is solvable. As $K(\zeta)/F$ is Galois and $G(K(\zeta)/F(\zeta, u_1))$ is solvable, the
group $G(K(\zeta)/F)$ is solvable by Remark 59.6. As $K/F$ is Galois, the group $G(K/F)$ is solvable again by Remark 59.6. □

The theorem shows that if $f$ is an irreducible polynomial in $F[t]$, e.g., $F = \mathbb{Q}$, then a formula for a root of $f$ exists that consists of addition, multiplication, and extraction of $n$th roots for various $n$, if and only if there exists a radical extension $K$ of $F$ such that $f$ has a root in $K$. As the conjugate of a radical extension is radical, if such a formula for one of the roots of irreducible $f_i$, such a formula holds for each root of $f_i$ in the appropriate conjugate.

Let $F$ be a field and $f$ be a non-constant polynomial in $F[t]$. We say that $f$ is solvable by radicals if for each irreducible factor $f_i$ of $f$ in $F[t]$, there exists a formula for a root of $f_i$ involving addition, multiplication, and extraction of $n$th roots for various $n$, i.e., there exists a radical extension $K_i$ of $F$ such that $f_i$ has a root in $K_i$ for each $i$.

We analyze what we have shown. Let $F$ be a field of characteristic zero. Suppose that $f$ in $F[t]$ is solvable by radicals and $f_i$ a irreducible factor of $f$ in $F[t]$. Then there exists a radical extension $K_i/F$ such that $f_i$ has a root in $K$. Assume that each such $K_i$ lies in a common extension field $\tilde{F}$, e.g., an algebraic closure of $F$. If there exist $m$ distinct (non-associative) factors $f_i$ of $f$, let $L = F(K_1 \cup \cdots \cup K_m)$, a radical extension of $F$. The normal closure $\tilde{L}/L$ in $\tilde{F}$ of $L/F$ is also a radical extension, so there exists a radical Galois extension $\tilde{L}/F$ such that $f$ splits over $\tilde{L}$. Let $\tilde{L}/E/F$ with $E$ a splitting field of $f$ over $F$. Although $E/F$ may not be radical, by the theorem, we know that $G(E/F)$ is solvable. Recall if $f$ is a polynomial in $F[t]$ and $K/F$ is a splitting field of $f$, then $G(K/F)$ is called the Galois group of $f$.

In general it is only unique up to isomorphism, but if we have fixed an algebraic closure of $F$, as we have done this time, it is unique. Therefore, we have shown the following:

**Theorem 59.9.** Let $F$ be a field of characteristic zero and $f$ an non-constant polynomial in $F[t]$. If $f$ is solvable by radicals, then the Galois group of $f$ is solvable.

We wish to show that there exist polynomials in $\mathbb{Q}[t]$ whose Galois groups are not solvable by radicals. We need the following group theory result and its application to field theory.

**Lemma 59.10.** Let $p$ be a prime and $G$ a subgroup of the symmetric group $S_p$ containing a $p$-cycle and a transposition. Then $G = S_p$.

This is not true if $p$ is not a prime.
To prove the lemma, we recall the following fact: For any \( n \)-cycle \((a_1 \cdots a_n)\), we have
\[
\sigma(a_1 \cdots a_n)\sigma^{-1} = (\sigma(a_1) \cdots \sigma(a_n))
\]
for all permutations \( \sigma \) in \( S_N \) with \( n \leq N \).

**Proof.** (of the lemma) By changing notation, we may assume that \((1 2) \in G\). Let \( \sigma = (a_1 \cdots a_p) \) lie in \( G \). As \( p \) is a prime, \( \sigma^i, 1 \leq i < p \), is also a \( p \)-cycle (cf. Exercise 24.24(1), so we may assume that \( \sigma = (1 2 \cdots p) \)). As \( \sigma(1 2)\sigma^{-1} = (2 3), \sigma(2 3)\sigma^{-1} = (3 4) \), etc., we see that \( (i i + 1) \) lies in \( G \) for all \( i \). Therefore \( (1 i + 1) = (1 i)(i i + 1)(1 i) \) lies in \( G \) for all \( i \) by induction, hence \( G = S_p \) by Proposition 24.7.

**Proposition 59.11.** Let \( f \) be an irreducible polynomial in \( \mathbb{Q}[t] \) of degree \( p \) with \( p \) a prime. Suppose that \( f \) has precisely two nonreal roots in the splitting field \( K \) in \( \mathbb{C} \) of \( f \) over \( \mathbb{Q} \). Then \( G(K/\mathbb{Q}) \cong S_p \).

**Proof.** Let \( \alpha \) in \( K \) be a root of \( f \), so \([\mathbb{Q}(\alpha) : \mathbb{Q}] = \deg f = p \). As \( \text{char } K = 0 \), we know that \( K/\mathbb{Q} \) is Galois, hence \( |G(K/\mathbb{Q})| = [K : \mathbb{Q}] \leq p! \). As \( G(K/\mathbb{Q}) \) acts as permutations on the roots of \( f \) and \( K \) is the splitting field of \( f \), we may assume that \( G(K/\mathbb{Q}) \subset S_p \). Since \( f \) has precisely two nonreal roots, the restriction of complex conjugation to \( K \) is a nontrivial (\( \mathbb{Q} \)-) automorphism of \( K \) and it must interchange these two nonreal roots and fix all the other roots of \( f \), i.e., it corresponds to a transposition in \( G(K/\mathbb{Q}) \). As \( p \mid [K : \mathbb{Q}] = |G(K/\mathbb{Q})| \), the group \( G(K/\mathbb{Q}) \) must contain an element of order \( p \) by Cauchy’s Theorem 21.23. But the only elements in \( S_p \) of order \( p \) are \( p \)-cycles. Therefore, \( G(K/\mathbb{Q}) = S_p \) by the lemma.

As a consequence, we obtain the following famous theorem:

**Theorem 59.12.** (Abel-Ruffini) In general, there is no formula to determine the roots of a general fifth degree polynomial in \( \mathbb{Q}[t] \) that involves only addition, multiplication, and extraction of \( n \)th roots for various \( n \).

**Proof.** It suffices to give an example of a polynomial \( f \) in \( \mathbb{Q}[t] \) that has a nonsolvable Galois group. Let \( f = t^5 - 6t + 3 \). We know that \( f \) is irreducible by Eisenstein’s Criterion, and by calculus it has three real roots in \( \mathbb{R} \), since \( f' = 5t^4 - 6 \) has only two real roots, \( f(-1) > 0 \), and \( f(1) < 0 \). By the proposition, the Galois group of \( f \) is isomorphic to \( S_5 \) which is not solvable.
Corollary 59.13. In general, there is no formula to determine the roots of a general $n$th degree polynomial in $\mathbb{Q}[t]$ that involves only addition, multiplication, and extraction of $n$th roots for various $n$.

**Proof.** Apply the theorem to any $n$th degree polynomial with $n \geq 5$ divisible by a fifth degree polynomial not solvable by radicals. □

Remark 59.14. 1. For each positive integer $n$, we have constructed fields $F$ and $K_n/F$ Galois with $G(K_n/F) \cong S_n$ in Example 56.6(3), so, in general, no formula can exist.

2. Hilbert showed that for each positive integer $n$, there exists a Galois extension $K_n/\mathbb{Q}$ with $G(K_n/\mathbb{Q}) \cong S_n$. This is not easy. It uses his Irreducibility Theorem. For a proof see Chapter 67. By the Primitive Element Theorem, $K_n = \mathbb{Q}(u)$ for some $u$, hence the Galois group of $m_\mathbb{Q}(u)$ is isomorphic to $S_n$, so is not solvable for $n \geq 5$.

3. There exist fields of characteristic zero such that every algebraic extension is radical, hence every non-constant polynomial over it is solvable by radicals, e.g., $\mathbb{R}$, $\mathbb{C}$. There are others, e.g., the so called local fields of characteristic zero that arise in number theory.

4. If $F$ is a finite field, then any finite extension of $F$ is cyclic so radical. In particular, the Galois group of any polynomial over $f$ is solvable by radicals.

We want a converse to our result. To do so needs additional concepts.

Definition 59.15. Let $K/F$ be a finite Galois extension of fields and $x$ an element of $K$. Define the norm of $x$ to be

$$N_{K/F}(x) := \prod_{G(K/F)} \sigma(x)$$

and the trace of $x$ to be

$$\text{Tr}_{K/F}(x) := \sum_{G(K/F)} \sigma(x).$$

Examples 59.16. Let $F$ be a field of characteristic different from two and $d$ an element in $F$ that is not a square. Let $K = F(\sqrt{d})$. If $x = a + b\sqrt{d}$ with $a, b$ in $F$, then $N_{K/F}(x) = a^2 - b^2 d$ and $\text{Tr}_{K/F}(x) = 2a$. Both of these lie in $F$. If, in addition, $x$ does not lie in $F$, i.e., $b$ is not zero, then $K = F(x)$ and $m_F(x) = t^2 - \text{Tr}_{K/F}(x)t + N_{K/F}(x) = (t - (a + b\sqrt{d}))(t - (a - b\sqrt{d}))$ lies in $F[t]$.

Properties 59.17. Let $K/F$ be a finite Galois extension of fields, $G = G(K/F)$, and $x, y$ elements of $K$. 

1. If \( \tau \) lies in \( G \), then \( \tau G = G = G \tau \), so
\[
\tau( N_{K/F}(x)) = N_{K/F}(x) = N_{K/F}(\tau(x))
\]
\[
\tau( Tr_{K/F}(x)) = Tr_{K/F}(x) = Tr_{K/F}(\tau(x)).
\]
In particular,
\[
N_{K/F}(x) \text{ and } Tr_{K/F}(x) \text{ lie in } K^G = F:
\]
We have
\[
\tau\left(\prod_G \sigma(x)\right) = \prod_G \tau \sigma(x) = \prod_G \sigma(x) = \prod_G \sigma \tau(x)
\]
\[
\tau\left(\sum_G \sigma(x)\right) = \sum_G \tau \sigma(x) = \sum_G \sigma(x) = \sum_G \sigma \tau(x).
\]

2. Suppose that \( G = \{\sigma_1, \ldots, \sigma_n\} \). As before let \( s_j \) denote the \( j \)th elementary symmetric function \( \sum_{1 \leq i_1 < \cdots < i_j \leq n} t_{i_1} \cdots t_{i_j} \) (and \( s_0 = 1 \)).
Then
\[
m_F(x) = \prod (t - \sigma_i(x)) = \sum (-1)^{n-i} s_{n-i}(\sigma_1(x), \ldots, \sigma_n(x)) t^i
\]
in \( F[t] \). In particular, if \( K = F(x) \), then
\[
N_{K/F}(x) = s_n(\sigma_1(x), \ldots, \sigma_n(x))
\]
\[
Tr_{K/F}(x) = s_1(\sigma_1(x), \ldots, \sigma_n(x)).
\]
What can you say if \( F(x) < K? \)

3. If \( x \) lies in \( F \), then \( \sigma(x) = x \) for each \( \sigma \) in \( G \). It follows that if \( K/E/F \) then
\[
N_{E/F}(x) = x^{[E:F]} \text{ and } Tr_{E/F}(x) = [E:F]x,
\]
as every \( F \)-embedding of \( E \) into \( K \) lifts to an element of \( G \).

4. The norm is multiplicative, i.e., \( N_{K/F}(xy) = N_{K/F}(x)N_{K/F}(y) \) and, if \( x \) is nonzero, \( N_{K/F}(x^{-1}) = (N_{K/F}(x))^{-1} \) (as \( N_{K/F}(1) = 1 \)). In particular,
\[
N_{K/F} : K^\times \to F^\times \text{ is a group homomorphism.}
\]

5. The map \( Tr_{K/F} : K \to F \) is an \( F \)-linear functional of \( F \)-vector spaces. By Dedekind’s Lemma \( 53.3 \), \( Tr_{K/F} \) is not the trivial map as it is the sum of distinct \( F \)-homomorphisms. In particular, \( \dim \ker Tr_{K/F} = [K : F] - 1 \).

We can generalize the norm and trace as follows:
Remark 59.18. Let $E/F$ be a finite separable extension of fields of degree $n$ with $K/E$ a normal closure of $E/F$. Hence $K/F$ is Galois. Let

$$
\tau_1, \ldots, \tau_m : E \to K
$$

be all the distinct $F$-homomorphisms. We know that $m = n$ as $E/F$ is finite and separable. Let $x$ be an element in $E$. Define the norm of $x$ to be

$$
N_{E/F}(x) := \prod_i \tau_i(x)
$$

and the trace of $x$ to be

$$
Tr_{E/F}(x) := \sum_i \tau_i(x).
$$

Since $K/F$ is Galois, if $\sigma \in G(K/F)$, then

$$
\sigma \tau_1, \ldots, \sigma \tau_n : E \to K
$$

are distinct $F$-homomorphisms, so these are just a permutation of $\tau_1, \ldots, \tau_n$, hence

$$
N_{E/F}(x) \text{ and } Tr_{E/F}(x) \text{ lie in } K^G = F
$$

just as in the Galois case. [Note we only needed $K/E/F$ with $K/F$ finite and Galois to deduce this i.e., is independent of the finite Galois extension of $F$ containing $E$.] Moreover, Properties 59.17 hold for $E/F$ with $\{\tau_1, \ldots, \tau_m\}$ replacing the $\sigma$ in Property 59.17(1) and $\{\sigma_1, \ldots, \sigma_n\}$ in Property 59.17(2) This generalization is useful in applications, especially in studying towers of finite separable extensions of fields. (Cf. the exercises at the end of this section). This can be further generalized to include inseparable extensions, but we shall not pursue this here.

In his book on algebraic number theory, Hilbert had the following result as his 90th Satz (theorem):

Lemma 59.19. (Hilbert Theorem 90) Let $K/F$ be a Galois extension of fields of degree $n$. Suppose that $K/F$ is cyclic with $G(K/F) = \langle \sigma \rangle$ and $x \in K$. Then

$$
N_{K/F}(x) = 1 \text{ if and only if there exists a } y \in K^\times \text{ satisfying } x = \frac{y}{\sigma(y)},
$$

i.e., we have an exact sequence

$$
1 \to K^\times \xrightarrow{\psi} K^\times \xrightarrow{N_{K/F}} F^\times
$$

with $\psi(y) = y/\sigma(y)$. 

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**Proof.** ($\Leftarrow$): If $x = y/\sigma(y)$, then

$$N_{K/F}(x) = \frac{N_{K/F}(y)}{N_{K/F}(\sigma(y))} = 1.$$  

($\Rightarrow$): Suppose that

$$1 = N_{K/F}(x) = x\sigma(x)\cdots\sigma^{n-1}(x).$$

In particular, $x$ is nonzero. By Dedekind’s Lemma 53.3, we know that the $F$-automorphisms $1, \sigma, \ldots, \sigma^{n-1}$ are independent, so if $a_0, \ldots, a_{n-1}$ lie in $K$, we have $\sum_{i=0}^{n-1} a_i \sigma^i(z) = 0$ for all $z$ in $K$ if and only if all the $a_i$ are zero. Set $a_0 = 1$ and $a_i = x\sigma(x)\cdots\sigma^{i-1}(x)$ in $K^\times$ for $i = 1, \ldots, n-1$. Then

$$g = 1 + a_1 \sigma + \cdots + a_{n-1} \sigma^{n-1}$$

$$= 1 + x\sigma + x\sigma(x)\sigma^2 + \cdots + x\sigma(x)\cdots\sigma^{n-2}(x)\sigma^{n-1}$$

is not trivial, so there exists an element $u$ in $K$ satisfying

$$0 \neq y := g(u) = u + x\sigma(u) + \cdots + x\sigma(x)\cdots\sigma^{n-2}(x)\sigma^{n-1}(u).$$

Since $\sigma^n = 1_K$ and $1 = N_{K/F}(x) = x\sigma(x)\cdots\sigma^{n-1}(x)$, we see that

$$x\sigma(y) = x\sigma(u) + x\sigma(x)\sigma^2(u) + \cdots + x\sigma(x)\cdots\sigma^{n-1}(x)\sigma^n(u) = y,$$

i.e., $x = y/\sigma(y)$ as needed. $\square$

Although the computational proof reveals little, generalizations of the lemma above have resulted in very deep theorems in algebra. Hilbert established it to prove the following important result.

**Theorem 59.20.** Let $K/F$ be a Galois extension of fields of degree $n$ with the characteristic of $F$ either zero or char $F \nmid n$. Suppose that $K/F$ is cyclic with $G(K/F) = \langle \sigma \rangle$ and $t^n - 1$ splits over $F$. Then there exists an element $u$ in $K$ satisfying $K = F(u)$ and $m_F(u) = t^n - a$ in $F[t]$ for some $a$ in $F$. In particular, $K/F$ is radical.

**Proof.** Let $\zeta$ be a primitive $n$th root of unity in $F$. By our previous work, we know that $U = \langle \zeta \rangle$ is a cyclic group satisfying $|U| = n$, so $|U| = [K : F]$. As $\zeta^{-1}$ lies in $F$ and $N_{K/F}(\zeta^{-1}) = \zeta^{-n} = 1$, an application of Hilbert Theorem 90 shows that there exists a nonzero element $u$ in $K$ satisfying $\zeta = \sigma(u)/u$, i.e., $\sigma(u) = \zeta u$. Since $\sigma(u^n) = \sigma(u)^n = \zeta^n u^n = u^n$, we have $\sigma^i(u^n) = u^n$ for all $i$. Therefore, $a := u^n$ lies in $K^{G(K/F)} = F$ as $K/F$ is Galois. Let $f := t^n - a$, a polynomial lying in $F[t]$. This polynomial has $n$ distinct roots: $\zeta^i u$ for $i = 1, \ldots, n$. In particular, $F(u)$ is a splitting field of $f$ over $F$. To finish, we need to show that $K = F(u)$, as then $\deg f = \deg m_F(u)$. We know that

$$\sigma^i|_{F(u)} : F(u) \to K$$

maps $u \mapsto \zeta^i u$
comprise \( n \) distinct \( F \)-homomorphisms, so by Artin’s Lemma 53.8, we have
\[
[F(u) : F] \geq n = |U| = [K : F],
\]
hence \( K = F(u) \).
\( \square \)

Let \( K/F \) be a Galois extension of fields of degree \( n \) with the characteristic of \( F \) either zero or \( \text{char} F \nmid n \). Suppose that \( K/F \) is cyclic with \( G(K/F) = \langle \sigma \rangle \) and \( t^n - 1 \) splits over \( F \). Then there exists an element \( u \) in \( K \) satisfying \( K = F(u) \) and \( m_F(u) = t^n - a \) in \( F[t] \) for some \( a \) in \( F \). In particular, \( K/F \) is radical.

**Corollary 59.21.** Let \( K/F \) be a cyclic extension of fields of degree \( n \) with the characteristic of \( F \) either zero or \( \text{char} F \nmid n \). Suppose that \( K/F \) is cyclic and \( t^n - 1 \) splits over \( F \). Then there exists an element \( x \in K \) such that \( \{ \sigma(x) \mid \sigma \in G(K/F) \} \) is an \( F \)-basis for \( K \). In particular, this is true if \( F \) is a finite field.

**Proof.** Let \( G(K/F) = \langle \sigma \rangle \). By the proof of Theorem 59.20, \( K = F(u) \) with \( u^i = \sigma(u^i) \) for \( i = 0, \ldots, n-1 \). The first statement follows. If \( F \) is a finite field, then \( L/F \) is a cyclic extension with \( \text{char} F \nmid |F| \) by Exercise 54.9(2).
\( \square \)

We shall show that this corollary generalizes to the case that \( K/F \) is a finite Galois extension in Theorem 60.5 below.

Theorem 59.20 allows us to deduce the desired appropriate converse of Theorem 59.8.

**Theorem 59.22.** Suppose that \( F \) is a field of characteristic zero and \( K/F \) is a finite Galois extension with \( G(K/F) \) solvable. Then there exists a finite extension \( L/K \) such that \( L/F \) is Galois and radical.

**Proof.** Let \( n = [K : F] = |G(K/F)| \).

**Case 1.** The polynomial \( t^n - 1 \) splits over \( F \):

We show in this case that \( K/F \) itself is radical and Galois. As the group \( G(K/F) \) is finite and solvable, it is polycyclic, i.e., has a cyclic series say
\[
1 = N_1 \lhd N_2 \lhd \cdots \lhd N_r = G
\]
with each \( N_i/N_{i-1} \) cyclic. Let \( F_i = K^{N_i} \). By the Fundamental Theorem of Galois Theory 56.3, we conclude that
\( K/F_{i+1} \) is Galois with \( G(K/F_{i+1}) = N_{i+1} \),
and, as \( N_i \lhd N_{i+1} \),
\( F_i/F_{i+1} \) is normal with \( G(F_i/F_{i+1}) \cong N_{i+1}/N_i \) cyclic.
Let \( n_i = |N_{i+1}/N_i| = [F_i : F_{i+1}] \). Since \([F_i : F_{i+1}] \mid [K : F] = |G(K/F)|\), we know that \( n_i \mid n\), hence \( t^{n_i} - 1 \mid t^n - 1 \) in \( F[t] \) (why?). Thus \( t^{n_i} - 1 \) in \( F[t] \subset F_{i+1}[t] \) splits over \( F_i \) for each \( i \). It follows that \( K/F \) is radical and Galois by Theorem 59.20.

**Case 2.** \( t^n - 1 \) in \( F[t] \) does not split over \( F \):

Let \( E \) be a splitting field of \( t^n - 1 \) over \( K \) and \( \zeta \) a primitive \( n \)th root of unity in \( E \).

**Claim:** \( K(\zeta)/F \) is a radical and Galois extension (and, therefore, we are done).

We have already seen that \( K(\zeta)/F \) is Galois as \( K/F \) is the splitting field of a separable polynomial over \( F \). So we need only show that \( K(\zeta)/F \) is a radical extension. As \( K(\zeta)/F(\zeta) \) is Galois and \( F(\zeta)/F \) is radical and Galois, it suffices to show that \( K(\zeta)/F(\zeta) \) is radical.

We reduce to Case 1 by showing:

(i) \([K(\zeta) : F(\zeta)] = |G(K/F(\zeta))| \mid n\).

(ii) \( G(K(\zeta)/F(\zeta)) \) is a solvable group.

Indeed, suppose we have shown (i) and (ii). Then \( F(\zeta) \) contains a primitive \(|G(K(\zeta)/F(\zeta))|\)th primitive root of unity, so by Case 1, the extension \( K(\zeta)/F(\zeta) \) is Galois and radical. As \( F(\zeta)/F \) is radical, it then follows that \( K(\zeta)/F \) is radical. So we need only show (i) and (ii).

Let \( \sigma \in G(K(\zeta)/F(\zeta)) \). Then \( \sigma \) fixes \( F \) (and \( \zeta \)), hence \( \sigma|_K : K \to K(\zeta) \) is an \( F \)-homomorphism. By Artin’s Lemma 53.8, there exist at most \([K : F] = |G(K/F)|\) \( F \)-homomorphisms \( K \to K(\zeta) \). As each element of \( G(K/F) \) is such an \( F \)-homomorphism, we must have \( \sigma|_K \) lies in \( G(K/F) \).

Therefore,

\[ \psi : G(K(\zeta)/F(\zeta)) \to G(K/F) \]

given by the restriction \( \sigma \mapsto \sigma|_K \)

is a well-defined group homomorphism. But \( \psi \) is monic, as \( \sigma|_K = 1_K \) means that \( \sigma \) fixes \( K \) and \( \zeta \) hence \( K(\zeta) \), i.e., \( \sigma = 1_{K(\zeta)} \). Therefore,

\[ |G(K(\zeta)/F(\zeta))| \mid |G(K/F)| = n \]

and \( G(K(\zeta)/F(\zeta)) \) is isomorphic to a subgroup of the solvable group \( G(K/F) \), so also solvable.

**Corollary 59.23.** Let \( F \) be a field of characteristic zero and \( f \) a non-constant polynomial in \( F[t] \). Then \( f \) is solvable by radicals if and only if the Galois group of \( f \) is solvable.

**Proof.** (\( \Rightarrow \)) has already been done.

(\( \Leftarrow \)): Let \( K \) be a splitting field of \( f \) over \( F \). By hypothesis, \( G(K/F) \) is solvable, so by the theorem there exists an extension \( L/K \) with \( L/F \) Galois and radical. Since \( f \) splits over \( L \), it is solvable by radicals. \( \square \)
Corollary 59.24. Let $F$ be a field of characteristic zero and $f$ a non-constant polynomial in $F[t]$ of degree at most four. Then $f$ is solvable by radicals.

Proof. The Galois group of $f$ is isomorphic to a subgroup of $S_4$, a solvable group. □

Remarks 59.25. The quadratic formula for the general polynomial of degree two is well-known and gives a root for any quadratic polynomial (and works over any field of characteristic not two). The corollary says if $F$ is a field of characteristic zero and $f$ the general polynomial of degree three (respectively, four), then there exists a formula involving only field operations and taking various $n$th roots giving a root, hence this formula is applicable to all polynomials of this degree. This is the formula of Ferro-Tartaglia-Cardano (respectively, Ferrari). Thus all the roots of polynomials of degree at most four over a field of characteristic zero can be found by such formulae.


1. Let $L/F$ be a field extension and $L/K_i/F$ with $K_i/F$ abelian. Show that $K_1 K_2 = K_1(K_2)$ is abelian.

2. Show that Lemma 59.10 is false if $p$ is not a prime.

3. Let $K/F$ be a finite separable extension of fields. Show that $N_{K/F} = N_{E/F} \circ N_{K/E}$ and $\text{Tr}_{K/F} = \text{Tr}_{E/F} \circ \text{Tr}_{K/E}$ for any intermediate field $K/E/F$.

4. Let $K/F$ be a finite separable extension of fields, $x$ an element of $K$, and $\lambda_x : F(x) \rightarrow F(x)$ the $F$-linear transformation given by $y \mapsto xy$. Show that the characteristic polynomial of $\lambda_x$ is $(m_F(x))_{[K:F(x)]}$. In particular, if $m_F(x) = t^r + a_{r-1} t^{r-1} + \cdots + a_0$ in $F[t]$, then $N_{K/F}(x) = (-1)^r a_0_{[K:F(x)]}$ and $\text{Tr}_{K/F}(x) = -[K : F(x)] a_{r-1}$.

5. Suppose that $K/F$ is a finite separable extension of fields, $K^* := \text{Hom}_F(K, F)$, the $F$-linear dual space of $K$. Show that the map $T : K \rightarrow K^*$ defined by $x \mapsto \text{tr}_x : y \mapsto \text{Tr}_{K/F}(xy)$ is an $F$-linear isomorphism. In particular, if $\{x_1, \ldots, x_n\}$ is an $F$-basis for $K$, then there exists a basis $\{y_1, \ldots, y_n\}$ for $K$ satisfying $\text{Tr}_{K/F}(x_i y_j) = \delta_{ij}$ (the Kronecker delta).

6. Let $K/F$ be a finite separable extension of fields, $x$ an element of $K$, and $\lambda_x : F(x) \rightarrow F(x)$ the $F$-linear transformation given by $y \mapsto xy$. If $B$ is an $F$-basis for $K$, show that $\text{Tr}_{K/F}(x) = \text{trace}[\lambda_x]_B$ and $N_{K/F} = \det[\lambda_x]_B$. 

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7. Let $K/F$ be a Galois extension of fields of degree $n$. Suppose that $K/F$ is cyclic with $G(K/F) = \langle \sigma \rangle$ and $x \in K$. Show $\text{Tr}_{K/F}(x) = 0$ if and only if there exists a $y \in K^\times$ satisfying $x = y - \sigma(y)$.

8. Let $F$ be a field of characteristic zero and $K/F$ an abelian extension. Let $n$ be a positive integer such that $\sigma^n = 1$ for every $\sigma \in G(K/F)$ and $F$ contains a primitive $n$th root of unity. Show that there exist $\alpha_1, \ldots, \alpha_r$ in $K$ such that $K = F(\alpha_1, \ldots, \alpha_r)$ with $\alpha_i^n \in F$ for some $n_i \mid n$ for $i = 1, \ldots, r$.

60. Addendum: On Hilbert Theorem 90

In this section we use the proof of Hilbert Theorem 90 (Theorem 59.19) to state a more modern version whose generalizations has become very important in algebra and number theory.

Let $K$ be a field and $G$ a group of automorphisms of $K$. We call $S_G := \{x_\sigma \in K \mid \sigma \in G\}$ a solution set to the Noether equations if $0 \not\in S_G$ and

\[ (*) \quad x_\sigma \sigma(x_\tau) = x_{\sigma \tau}, \text{ for all } \sigma, \tau \in G. \]

If $x_\sigma = 0$ satisfies $(*)$, then $S_G = \{0\}$, since $\tau G = G$ for all $\tau \in G$. This is the reason that we exclude $0 \in S_G$.

**Theorem 60.1.** Let $K$ be a field and $G$ a finite group of automorphisms of $K$. Then a set $S = \{x_\sigma \in K \mid \sigma \in G\}$ is a solution set to the Noether equations if and only if there exists an element $y$ in $K$ satisfying $x_\sigma = y/\sigma(y)$ for all $\sigma \in G$.

The proof is almost identical to the proof of Hilbert 90 59.19.

**Proof.** If $y$ in $K^\times$ satisfies $x_\sigma = y/\sigma(y)$ for all $\sigma \in G$, then

\[
\frac{y}{\sigma(y)} \sigma \left( \frac{y}{\tau(y)} \right) = \frac{y}{\sigma(y)} \frac{\sigma(y)}{\sigma \tau(y)} = \frac{y}{\sigma \tau(y)}
\]

for all $\sigma, \tau \in G$. So $S = \{y/\sigma y \mid \sigma \in G\}$ is a solution set to the Noether equations.

Conversely, suppose that $S$ is a solution set to the Noether equations. By Dedekind’s Lemma 53.3, $G$ is not an independent set, so there exists an element $y \in K^\times$ satisfying $0 \neq z = \sum_{\tau \in G} x_\tau \tau(y)$. Hence $\sigma(z) = \sum_{\tau \in G} \sigma(x_\tau) \sigma \tau(y)$. Multiply this equation by $x_\sigma$. Using $(*)$ then yields

\[ x_\sigma \sigma(z) = \sum_{\tau \in G} x_\sigma \sigma(x_\tau) \sigma \tau(y) = \sum_{\tau \in G} x_\tau \tau(y) = y. \]

Therefore, $x_\sigma = y/\sigma(y)$ for all $\sigma \in G$. \qed
We now construct the general setup where the Noether equations constitute a special case.

**Construction 60.2.** Let $G$ be a multiplicative group, $A$ a (multiplicative) abelian group, and $G \to \text{Aut}(A)$ a group homomorphism. This defines a $G$-action on $A$ by evaluation that we write by $\sigma(a)$ for $\sigma \in G$ and $a \in A$. We call a map $f : G \to A$ a *crossed homomorphism*, if

$$f(\sigma \tau) = f(\sigma) \sigma(f(\tau))$$

for all $\sigma, \tau \in G$. In particular, if $G \to \text{Aut}(A)$ is the trivial map, the set of crossed homomorphisms is just $\text{Hom}(G, A)$, the group of group homomorphism $G \to A$. A crossed homomorphism $f : G \to A$ is called a *principal crossed homomorphism* if there exists a $y \in A$ satisfying $f(\sigma) = \sigma(y)/y$ for all $\sigma \in G$.

Set

$$Z^1(G, A) := \{ f : G \to A \mid f \text{ a crossed homomorphism} \}$$

$$B^1(G, A) := \{ f : G \to A \mid f \text{ a principal crossed homomorphism} \}.$$

Then $Z^1(G, A)$ is an abelian group and $B^1(G, A)$ is a subgroup (whose proof we leave as an exercise). The quotient group

$$H^1(G, A) = Z^1(G, A)/B^1(G, A)$$

is called the 1st cohomomology group of $G$ with coefficients in $A$.

**Corollary 60.3.** (Hilbert Theorem 90) Let $L/F$ be a finite Galois extension. Then

$$H^1(G(L/F), L^\times) = 1.$$

**Proof.** This is just a reformulation of Theorem 60.1. □

It can be shown, if $L/F$ is any Galois extension, that

$$H^1_{\text{cont}}((G(L/F), L^\times) = 1,$$

where all the crossed homomorphisms defined above must also be continuous relative to the topology introduced in Section 57. In particular, if $L/F$ is finite, $G(L/F)$ has the discrete topology and this agrees with the corollary.

The analogue for Hilbert Theorem 90 for the additive group of a finite cyclic Galois extension is also true with the trace replacing the norm. (Cf. Exercise 59.26(7).) The cohomological generalization is also true. Changing the notation above from an abelian multiplicative group to that of an additive group with the appropriate changes of notation that this forces results in the following:
Corollary 60.4. Let \( L/F \) be a finite Galois extension. Then
\[
H^1(G(L/F), L^+) = 0,
\]
where \( L^+ \) is the additive group of \( L \).

We leave its proof as an exercise.

Actually more is true. Not only is \( H^1(G(L/F), L^+) = 0 \), but all the higher cohomology groups \( H^n(G(L/F), L^+) \) (that we have not defined) are trivial for \( n \geq 1 \) for any finite Galois extension \( L/F \) (and all Galois extensions if use the continuous crossed homomorphisms). This is certainly not true for \( L \times \). [We shall investigate \( H^2(G(L/F), L^+) \) in Section 104.] The underlying mathematical reason for this is the following result that also generalizes a cyclic case, Corollary 59.21.

Theorem 60.5. (Normal Basis Theorem) Let \( L/F \) be a finite Galois extension. Then \( L \) has a normal basis, i.e., there exists an element \( x \) in \( L \) satisfying \( \mathcal{N} = \{ \sigma(x) \mid \sigma \in G(L/F) \} \) is a basis for \( L \) over \( F \).

Proof. Let \( L/F \) be of degree \( n \) and \( G(L/F) = \{ 1 = \sigma_1, \ldots, \sigma_n \} \).

If \( F \) is a finite field, the result follows by Corollary 59.21, so we may assume that \( F \) is an infinite field. In particular, we may identify polynomials over \( L \) and their associated polynomial functions.

Let \( t_i = t_{\sigma_i}, i = 1, \ldots, n \), be \( n \) distinct variables. Set \( t_{\sigma_i, \sigma_j} = t_{\sigma_i^{-1}\sigma_j} \) for \( i, j = 1, \ldots, n \). Let \( f \in L[t_1, \ldots, t_n] \) be the polynomial defined by
\[
f(t_1, \ldots, t_n) = \det (t_{\sigma_i, \sigma_j}).
\]
As \( L \) is infinite and the evaluation \( e_{(1,0,\ldots,0)}f = f(1,0,\ldots,0) \neq 0 \), the polynomial \( f \neq 0 \). Hence the polynomial function \( g : L[t] \rightarrow L \) defined by \( g(x) = f(\sigma_1(x), \ldots, \sigma_n(x)) \) cannot have infinitely many roots. It follows that there exists an \( x \) in \( L \) satisfying
\[
0 \neq f(\sigma_1(x), \ldots, \sigma_n(x)) = \det (\sigma_i^{-1}\sigma_j(x)).
\]
Suppose that there exist \( a_1, \ldots, a_n \) in \( F \) satisfying
\[
(*) \quad 0 = a_1\sigma_1(x) + \cdots + a_n\sigma_n(x).
\]
Applying \( \sigma_i^{-1} \) to \( (*) \), for \( i = 1, \ldots, n \), yields a system of equations:
\[
0 = a_1\sigma_i^{-1}\sigma_1(x) + \cdots + a_n\sigma_i^{-1}\sigma_n(x), \quad i = 1, \ldots, n.
\]
Since \( \det (\sigma_i^{-1}\sigma_j(x)) \neq 0 \), Cramer’s Rule implies that \( a_i = 0 \) for \( i = 1, \ldots, n \). It follows that \( \mathcal{N} \) is an \( F \)-basis for \( L \). \( \square \)

Exercises 60.6.
1. Let $L/F$ be a finite Galois extension. Let $y \in K$ satisfy $x_\sigma = y/\sigma(y)$ for all $x \in L^\times$. Let $r$ be the least common multiple of the orders of the elements of $G(L/F)$. The $y^r$ lies in $F$.

2. Let $G$ be a group and $A$ a (multiplicative) abelian group. Show that $Z^1(G, A)$ and $B^1(G, A)$ are abelian groups.

3. Let $L/F$ be a finite Galois extension and suppose that the group $\mu_n$ of $n$th roots of unity lies in $F$ and consists of $n$ elements. Then show that $B^1(G(L/F), \mu_n) = 1$ and $Z^1(G(L/F), \mu_n) = \text{Hom}(G(L/F), \mu_n)$, i.e., $H^1(G(L/F), L^\times) = \text{Hom}(G(L/F), \mu_n)$.

4. Prove Corollary 60.4.

### 61. Addendum: Kummer Theory

In Section 59, we proved Theorem 59.20 classifying finite cyclic field extensions, i.e., finite Galois extensions with cyclic Galois group, of degree $n$ when the base field contained all the $n$th roots of unity and the characteristic of the field was zero or did not divide $n$. Exercise 59.26(8) extends this to finite abelian extension of a field. In this section, we generalize this to arbitrary abelian extensions of a field, i.e., not necessarily finite. To do this, we shall need results from Section 57.

We begin with a preliminary discussion of pairing of groups. Let $A$, $B$, and $C$ be abelian multiplicative groups. A pairing $\langle \cdot, \cdot \rangle : A \times B \to C$ is a map that is a group homomorphism in each variable. Set

$$A^\perp := \{ a \in A \mid \langle a, x \rangle = e_C \text{ for all } x \in B \}$$

$$B^\perp := \{ b \in B \mid \langle y, b \rangle = e_C \text{ for all } y \in A \}.$$

This pairing induces a homomorphism $A \to \text{Hom}(B, C)$ by $a \mapsto \langle a, \cdot \rangle$, hence a monomorphism $A/A^\perp \to \text{Hom}(B/B^\perp, C)$. Similarly, this pairing induces a homomorphism $B \to \text{Hom}(A, C)$ by $b \mapsto \langle \cdot, b \rangle$, hence a monomorphism $B/B^\perp \to \text{Hom}(A/A^\perp, C)$.

Suppose $m$ is a positive integer such that both $A/A^\perp$ and $B/B^\perp$ satisfy $x^m = e_A$ and $y^m = e_B$ for all $x \in A$ and for all $y \in B$, and $C = \mu_m$ is the cyclic group of order $m$, i.e., $m$ divides the exponent of these three groups, where the exponent $N$ of a (multiplicative) group $G$ is the least positive integer $N$ such that $g^N = e_G$ for all $g \in G$ (if such exists). Then in the above,

$$A/A^\perp \text{ is finite if and only if } B/B^\perp \text{ is finite},$$

and if this is the case, then we have an isomorphism:

$$A/A^\perp \xrightarrow{\sim} B/B^\perp := \text{Hom}(B/B^\perp, \mu_m).$$
We shall apply this when $\mu_m$ is the group of $m$th roots of unity in a field $F$ of characteristic zero or characteristic not dividing $m$.

We set up notation to be used in the rest of this section. Let $F$ be a field. Denote by $\tilde{F}$ a fixed algebraic closure of $F$ and assume all algebraic extensions of $F$ lie in it. So the group of $m$th roots of unity $\mu_m$ lie $\tilde{F}$. We shall assume that $|\mu_m| = m$, so the characteristic of $F$ is zero or does not divide $m$. We shall also assume that $\mu_m \subset F$. Therefore, if $K/F$ is a splitting field of $t^m - a \in F[t]$, then $K = F(\alpha)$ with $\alpha^m = a$ and the $m$ distinct roots of $t^m - a$ are given by $\zeta_\alpha$ for $\zeta \in \mu_m$ by Theorem 59.20. We write this symbolically as $K = F(\{a_1^{1/m}\})$. Let $(F^\times)^m = \{x^m \mid x \in F\}$ and if $(F^\times)^m \subset P \subset F^\times$ is a subgroup, let $F_P = F(P^{1/m})$ be the field generated by the $F(a_1^{1/m})$ with $a \in P$. The (possibly infinite) field extension $F_P/F$ is normal and separable, hence is Galois by Theorem 57.3.

**Definition 61.3.** Let $K/F$ be a (possibly infinite) Galois extension of fields. We say that $K/F$ has exponent $m$ if $G(K/F)$ has finite exponent dividing $m$. Suppose, in addition, that $\mu_m \subset F$ and $|\mu_m| = m$. We say that $K/F$ is a Kummer extension of exponent $m$ if $K/F$ is abelian of exponent $m$, i.e., Galois with abelian Galois group.

Let $F$ be a field containing the $m$th roots of unity $\mu_m$ with $|\mu_m| = m$. Set

$$F_m := \{K \mid K/F \text{ is a Kummer extension of exponent } m\}$$

$$G_m^{ab} := \{P \mid (F^\times)^m \subset P \subset F^\times \text{ is a subgroup}\}.$$

**Theorem 61.4.** Let $F$ be a field containing the $m$th roots of unity $\mu_m$ with $|\mu_m| = m$. Then the map $i : G_m \to F_m$ given by $P \mapsto F_P$ is a bijection. Moreover, $F_P/F$ is a finite extension if and only if $P/(F^\times)^m$ is a finite group, and if finite, then $[F : F_P] = [P : (F^\times)^m]$.

**Proof.** Let $P \in G_m$. We first show that $F_P \in F_m$. Let $a \in P$, so $a = \alpha^m$, some $\alpha \in F$. In the proof of Theorem 59.20 using Hilbert Theorem 90 (Lemma 59.19), we saw that if $\sigma \in G(F_P/F)$, then there exists a unique $\zeta_\sigma \in \mu_m$ satisfying $\sigma(\alpha) = \zeta_\sigma \alpha$. Check that $\zeta_\sigma$ is independent of the choice of $\alpha$. In particular, $\sigma^m(\alpha) = \alpha$ and if $\tau \in G(F_P/F)$, then $\sigma\tau(\alpha) = \zeta_\sigma \zeta_\tau = \tau \sigma(\alpha)$.

It follows that $F_P/F$ is an abelian extension. Since $F_P = F(\{\alpha \in F \mid \alpha^m \in P\}$, we have $F_P \in F_m$.

Next let $K \in F_m$. Set

$$N := \{a \in F^\times \mid \text{there exists } \alpha \in K \text{ such that } \alpha^m = a\} \in G_m.$$
Let \( G = G(K/F) \). Define
\[
 f : N \rightarrow \hat{G} := \text{Hom}(G, \mu_m)
\]
as follows: Let \( a \in N \). Choose \( \alpha \in K \) satisfying \( \alpha^m = a \). If \( \sigma \in G \), then, as above, there exists a unique \( \zeta_\sigma \in \mu_m \) satisfying \( \sigma(\alpha) = \zeta_\sigma \alpha \).

Define \( f_a : N \rightarrow \hat{G} = \text{Hom}(G, \mu_m) \) by \( \sigma \mapsto \sigma(\alpha)/\alpha = \zeta_\sigma \).

**Claim 61.5.** \( f_a \) is a well-defined homomorphism.

If \( f_a \) is well-defined, it is clearly a homomorphism. So we must show that \( f_a \) is independent of the choice of \( \alpha \). We leave its proof as an exercise. By the claim we obtain a map \( f : N \rightarrow \hat{G} \) given by \( a \mapsto f_a \).

We apply this construction to \( K = F_P \) with \( P \in \mathcal{G}_m \). So we have a pairing
\[
\langle \ , \ \rangle : G \times P \rightarrow \mu_m
\]
defined by \( \langle \sigma, a \rangle = f_a(\sigma) = \sigma(\alpha)/\alpha \) where \( a = \alpha^m \).

**Claim 61.6.** \( G^\perp = 1 \) and \( P^\perp = (F^\times)^m \).

That \( G^\perp = 1 \) is easy, so we turn to \( P^\perp \). Suppose that \( \alpha \in K = F_P \) satisfies \( \alpha^m = \alpha \) and \( \langle \sigma, a \rangle = \sigma(\alpha)/\alpha = 1 \) for all \( \sigma \in G \), but \( \alpha \notin F \). Then there exists a \( \tau \in G \) such that \( \tau(\alpha) \neq \alpha \). Hence \( \langle \tau, a \rangle = \tau(\alpha)/\alpha \neq 1 \), a contradiction. The Claim now follows easily.

We also conclude by (61.1) that \( G \) is finite, i.e., \( F_P/F \) is finite, if and only if \( P/(F^\times)^m \) is finite. If this is the case then by (61.2), we have
\[(61.7) \quad [F : F_P] = [P : (F^\times)^m] \quad \text{and} \quad P/(F^\times)^m/(F^\times)^m \cong \hat{G}(F_P/F).
\]
This is the crucial observation.

**Claim 61.8.** \( P_1 \subset P_2 \) in \( \mathcal{G}_m \) if and only if \( F_{P_1} \subset F_{P_2} \). In particular, \( i : \mathcal{G}_m \rightarrow \mathcal{F}_m \) is injective.

The only if statement is clear, so suppose that \( F_{P_1} \subset F_{P_2} \). Let \( \alpha \in P_1 \). Then \( \alpha \) lies in \( F_{P_2} \), hence there exists a finitely generated extension \( K/F(F^m) \) with \( F_{P_2}/K \) and \( \alpha \in K \). It follows that we may assume that \( F_{P_2}/F(F^m) \) is finite generated, hence a finite field extension. Let \( P_3 \) be the group generated by \( \alpha \) and \( P_2 \). Then \( F(P_2) = F(P_3) \), and \( [F_{P_2} : F^m] = [F_{P_3} : F^m] \) by (61.7). Hence \( P_2 = P_3 \). This proves the claim.

So we are left to show that \( i : \mathcal{G}_m \rightarrow \mathcal{F}_m \) is surjective. Let \( K \in \mathcal{F}_m \), say \( K = F(P) \) with \( P \in \mathcal{G}_m \). For any element \( x \) in \( K \), there exists a finite subgroup \( P_0 \in \mathcal{G}_m \) such that \( x \) lies in \( F(P_0) \), i.e., \( K \) is the compositum of finite abelian extensions. Therefore, we may assume that \( K/F \) is finite abelian extension with \( G = G(K/F) \) of finite exponent dividing
Since $G$ is a finite abelian group of finite exponent dividing $m$, it is a product of finitely many cyclic groups of exponent dividing $m$ by Theorem 43.6 (or Proposition 14.9). Consequently, $G \cong \hat{G}$ by Exercise 43.24(10). By Exercise 56.13(13), this reduces us to the case that $K/F$ is finite cyclic which has been done. \hfill \Box

**Remark 61.9.** Let $F$ contain $\mu_m$ with $|\mu_m| = m$ as in the theorem. In the notation of the theorem, if we let $P \in \mathcal{G}_m$ and $G = G(F_P/F)$, then the map

\[ f : P \to \hat{G} \] given by \( a \mapsto (f_a : \sigma \to \sigma(\alpha)/\alpha) \) with $\alpha^m = a$

induces a monomorphism $\tilde{f} : F/(F^\times)^m \to \hat{G}$. If we take continuous homomorphisms $G \to \mu_m$ in the Krull topology on $G(\tilde{F}/F)$, then this map is an isomorphism. [As $\mu_m$ lies in $F$, its topology is discrete.] We leave this as an exercise. Write this as

\[
\begin{align*}
1 \to (F^\times)^m & \to P \to \text{Hom}_{\text{cont}}(G, \mu_m) \to 1 \\
\end{align*}
\]

is exact where $\text{Hom}_{\text{cont}}(G, \mu_m)$ denotes the group of continuous maps $G \to \mu_m$.

More generally, let $F_{\text{sep}}$ denote the *separable closure of* $F$ in $\tilde{F}$, i.e., the maximal separable extension of $F$ in $\tilde{F}$. Then $\Gamma_F := G(F_{\text{sep}}/F)$ is $G(\tilde{F}/F)$ and called the *absolute Galois group of* $F$. If $\mu_m \subset F$ and has $m$ elements, then we have an exact sequence

\[ 1 \to \mu_m \to F_{\text{sep}}^\times \overset{g}{\to} F_{\text{sep}}^\times \to 1. \]

where $g(x) = x^m$ By using ‘group cohomology’ — a subject that we shall not discuss — one obtains an exact sequence

\[ 1 \to \mu_m \to F^\times \overset{\delta}{\to} (F^\times) \overset{\delta}{\to} \text{Hom}_{\text{cont}}(\Gamma_F, \mu_m) \to 1. \]

called the *Kummer sequence* with $\delta$ a map coming from group cohomology theory called a connecting map. The surjectivity of $\delta$ arises from a homological version of Hilbert Theorem 90, the finite case can be found in Corollary 60.3.

**Exercises 61.12.**

1. Prove 61.1
2. Prove that $\zeta_\sigma$ in the proof of the theorem is independent of the choice of $\alpha$.
3. Prove Claim 61.5.
4. Prove that (61.10) implies (61.11).
62. Addendum: Galois’ Theorem

In this section, we establish Galois’ characterization of the Galois group of an irreducible polynomial of prime degree over a field of characteristic zero. We begin with a lemma, previously given as Exercise 22.15 (15). If $p$ is a prime, then a finite group is called an elementary $p$-group if it is isomorphic to $(\mathbb{Z}/p\mathbb{Z})^n$ for positive integer $n$. Note that any elementary $p$-group can be viewed as a finite dimensional vector space over $\mathbb{Z}/p\mathbb{Z}$. A minimal normal subgroup of a nontrivial group $G$ is a normal subgroup of $G$ such that the trivial group is the only normal subgroup of $G$ properly contained in it.

Lemma 62.1. Let $G$ be a nontrivial finite solvable group. Then any minimal normal subgroup of $G$ is an elementary $p$-group.

Proof. Let $H$ be a minimal normal subgroup of $G$. As $H$ is a subgroup of a solvable group, it is solvable. We also know that any characteristic subgroup of $H$ is normal in $G$ (cf. Exercise 11.9(20)). Since the commutator subgroup of any group is a characteristic subgroup, the series obtained by successively taking commutators gives a characteristic series for $H$ (cf. Proposition 17.5). As $H$ is solvable, it contains a characteristic nontrivial abelian normal subgroup, so we may assume $H$ is abelian. Since Sylow groups of an abelian group are characteristic, we may assume that $H$ is a $p$-group. Since $\{x \in H \mid x^p = 1\}$ is easily seen to be a characteristic subgroup of $H$, the lemma follows. □

More generally, it can be shown that a minimal normal subgroup of a nontrivial finite group is either simple or a product of simple groups all which are isomorphic.

Let $p$ be a prime and $V$ a finite dimensional vector space over $\mathbb{Z}/p\mathbb{Z}$. If $v_0$ is a vector in $V$, then the map $\tau_{v_0} : V \to V$ given by $v \mapsto v + v_0$ is called translation by $v_0$. [Note that it is not an $F$-linear transformation unless $v_0 = 0$.] The set of all translations of $V$ form a subgroup of the the group of affine transformations of $V$, the group defined by

$$\text{Aff}(V) := \{\tau : V \to V \mid \tau : v \mapsto av + v_0 \text{ with } a \text{ in } \mathbb{Z}/p\mathbb{Z} \text{ and } v_0 \text{ in } V\}.$$ 

So every element in $\text{Aff}(V)$ is a composition $\tau \lambda_a$ with $\tau$ a translation and $\lambda_a : V \to V$ the $F$-linear map given by $x \mapsto ax$ for some $a$ in $F$.

Theorem 62.2. Let $p$ be a prime and $S$ a set with $p$ elements. Suppose that $G$ is a transitive subgroup of $\Sigma(S) \cong S_p$. Then the following are equivalent:

1. $G$ is solvable.
(2) Each non-identity element in $G$ fixes at most one element in $S$.

(3) There exists a transitive normal subgroup $T$ of $G$ satisfying $|T| = p$ and $T$ is its own centralizer in $G$, i.e., $T = Z_G(T)$.

(4) $G$ is isomorphic to a subgroup of the group of affine transformations $\text{Aff}(\mathbb{Z}/p\mathbb{Z})$ and contains the subgroup of translations.

Moreover, if $G$ is solvable and if $T$ is as in (3), then there exists a cyclic subgroup $H$ in $G$ satisfying $G = HT$, $H \cap T = 1$, and $|H| | p - 1$. Furthermore, any non-identity normal subgroup of $G$ is of the form $H'T'$ for some subgroup $H'$ of $H$.

**Proof.** We may assume that $G \subset S_p$, so $S = \{1, \ldots, p\}$. As $G$ is transitive, the orbit $Gs$ for any $s \in S$, has order $p$, so $p | |G|$ and all the Sylow $p$-groups of $G$ are cyclic of order $p$, i.e., is generated by a $p$-cycle.

(1) $\Rightarrow$ (3): Let $N$ be a normal subgroup of $G$.

**Claim.** The group $N$ acts transitively on $S_p$. In particular, $p | |N|:

Let $x$ and $y$ be elements in $S$. As $G$ acts transitively on $S$, there exists an element $\sigma$ in $G$ satisfying $y = \sigma(x)$. As the map $Nx \to Ny$ given by $zx \mapsto z\sigma(x)$ is a well-defined bijection, we see that all orbits of $N$ have the same number of elements. If $O$ is a system of representatives for the action of $N$ on $S$, then $p = |S| = \sum_O |Nx| = \sum_O |N : N_x|$, by the Orbit Decomposition Theorem 19.9. So either $S$ has one orbit under the action of $N$, i.e., $N$ is transitive, or every orbit of $S$ under the action of $N$ has one point, i.e., $N$ is the identity subgroup. This establishes the claim.

Let $T$ be a minimal normal subgroup of $G$. By the lemma and claim, $T$ is a cyclic subgroup of order $p$, hence is generated by a $p$-cycle, say $T = \langle \sigma \rangle$. The normality of $T$ implies that it is the unique Sylow $p$-subgroup of $G$. Suppose $\tau$ centralizes $T$, i.e., lies in $Z_G(T)$. If $\tau(s) = s$ for some $s \in S$, then $\tau\sigma(s) = \sigma\tau(s) = \sigma(s)$, hence $\tau\sigma^i(s) = \sigma^i\tau(s) = \sigma^i(s)$ for all $i$ implying $\tau = 1$. Therefore $\tau$ is either 1 or is fixed point-free, so a $p$-cycle.

(3) $\Rightarrow$ (4): Suppose that $T = \langle \sigma \rangle$ is a cyclic subgroup of $G$ of order $p$.

We may assume that $S = \mathbb{Z}/p\mathbb{Z}$ and $\sigma$ is the $p$-cycle $(0 1 \cdots p - 1)$ where $- : \mathbb{Z} \to \mathbb{Z}/p\mathbb{Z}$ is the canonical epimorphism. Therefore, $\sigma$ generates the group of translations in $\text{Aff}(\mathbb{Z}/p\mathbb{Z})$. Let $\tau$ be an element in $G$. Then there exists an integer $n$ with $p \nmid n$, unique modulo $p$, such $\tau\sigma\tau^{-1} = \sigma^n$, i.e., we have a well-defined map $\varphi : G \to (\mathbb{Z}/p\mathbb{Z})^\times$ given by $\tau \mapsto n$. Clearly, this map is a group homomorphism with $\ker \varphi = Z_G(T)$. As
\[ \tau \sigma = \sigma^{\varphi(\tau)} \tau, \]  
we have 
\[ \tau(x + 1) = \tau \sigma(x) = \sigma^{\varphi(\tau)} \tau(x) = \tau(x) + \varphi(\tau). \]

In particular, \( \tau(1) = \tau(0) + \varphi(\tau) = \varphi(\tau) + b \) with \( b = \tau(0) \). Inductively, 
\[ \tau(x) = \varphi(\tau) x + b. \]
This yields (4).

(4) \( \Rightarrow \) (1): The translations form a cyclic normal subgroup \( T \) in \( G \) with quotient embedding in \( (\mathbb{Z}/p\mathbb{Z})^\times \).

(4) \( \Rightarrow \) (2): The equation 
\[ ax + b \equiv x \mod p, \]
with \( a \) and \( b \) integers, has a unique solution modulo \( p \) if either \( a \not\equiv 1 \mod p \) or \( b \not\equiv 0 \mod p \). It follows that any non-identity element in \( G \) can fix at most one point.

(2) \( \Rightarrow \) (3): We apply Exercise 21.26(17) — known as the Burnside Counting Theorem — to see that 
\[ |G| = \sum_{\tau \in G} |F_\tau(S)| \text{ where } F_\tau(S) := F_{(\tau)}(S). \]
As \( \tau = 1 \) fixes \( p \) points, we have \( \sum_{\tau \neq 1} |F_\tau(S)| = |G| - p \). By hypothesis, if \( \tau \neq 1 \) then \( |F_\tau(S)| \leq 1 \), so we must have precisely \( p - 1 \) non-identity elements \( \tau \) in \( G \) satisfying \( |F_\tau(S)| = 0 \), i.e., \( p - 1 \) elements in \( G \) have no fixed points in \( S \). If \( \sigma \) is one such, then so is \( \sigma^i \) for \( i = 1, \ldots, p - 1 \). It follows that \( \sigma \) is a \( p \)-cycle and \( T = \langle \sigma \rangle \), a group of order \( p \), contains all the \( p \)-cycles in \( G \). As the conjugate of a \( p \)-cycle is a \( p \)-cycle, \( T \triangleleft G \). Suppose that \( \tau \) lies in \( Z_G(T) \). If \( F_\tau(S) \) is empty, then \( \tau \) lies in \( T \), so assume that \( F_\tau(S) \) is not empty, say \( \tau(x) = x \). It follows that \( \tau \sigma(x) = \sigma \tau(x) = \sigma(x) \), i.e., both \( x \) and \( \sigma(x) \) lie in \( F_\tau(S) \). By assumption this means \( \tau = 1 \), so \( T = Z_G(T) \), establishing (3).

For the last two statements, let \( G \) satisfy (4). Set 
\[ H = \{ \tau \in G \mid \tau(x) \equiv ax \mod p \text{ with } a \text{ an integer satisfying } p \mid a \}. \]
Clearly, \( H \cap T = 1 \) and \( HT = G \). Certainly, we have a monomorphism \( H \to (\mathbb{Z}/p\mathbb{Z})^\times \), so \( H \) is cyclic and \( |H| = p - 1 \). If \( N \triangleleft G \) is not the trivial group, then by the claim, \( N \) acts transitively on \( S \) and contains \( T \), the unique Sylow \( p \)-subgroup of \( G \). It follows that \( N = H'T \) for some subgroup \( H' \) of \( H \) by the Correspondence Principle.

**Theorem 62.3.** (Galois) Suppose that \( F \) is a field of characteristic zero and \( f \) an irreducible polynomial in \( F[t] \) of prime degree \( p \). Let \( K \) be a splitting field of \( f \) over \( F \). Then \( f \) is solvable by radicals if and only if \( K = F(\alpha, \beta) \) for any two roots of \( f \) in \( K \). Moreover, if \( f \) is solvable by radicals, \( G(K/F) \) is isomorphic to a subgroup of
Aff($\mathbb{Z}/p\mathbb{Z}$) containing the subgroup of translations, and there exists an intermediate field $K/E/F$ with $E/F$ cyclic, $[K:E]=p$. Furthermore, if $K/L/F$ is an intermediate field with $L/F$ normal with $F<L$, then $L \subset E$.

**Proof.** We know that $G(K/F)$ can be viewed as a transitive subgroup of $S_p$ and $f$ is solvable by radicals if and only if $G(K/F)$ is solvable. By the previous theorem, $G(K/F)$ is solvable if and only if no non-identity element in $G(K/F)$ fixes two roots of $f$. Let $\alpha$ and $\beta$ be two roots of $f$ in $K$. By the Fundamental Theorem of Galois Theory 56.3, we know that $K=F(\alpha,\beta)$ if and only if $G(K/F(\alpha,\beta))=1$ if and only if no non-identity element of $G(K/F)$ fixes $\alpha$ and $\beta$. The result now easily follows. □

### 63. The Discriminant of a Polynomial

Let $R$ be a commutative ring and let the symmetric group $S_n$ act on $R[t_1,\ldots,t_n]$ by $\sigma t_i = t_{\sigma(i)}$ for all $i$ and $\sigma \in S_n$. This induces a group monomorphism $S_n \to \text{Aut}(R[t_1,\ldots,t_n])$ by viewing $\sigma \in S_n$ as a ring homomorphism fixing $R$. Analogous to the field case of $F(t_1,\ldots,t_n)^{S_n}$, one can show that the set of elements $R[t_1,\ldots,t_n]^{S_n}$ in $R[t_1,\ldots,t_n]$ fixed by $S_n$, is equal to $R[s_1(t_1,\ldots,t_n),\ldots,s_n(t_1,\ldots,t_n)]$, where $s_j(t_1,\ldots,t_n) = \sum_{1 \leq i_1 < \ldots < i_j \leq n} t_{i_1} \cdots t_{i_j}$ with $1 \leq j \leq n$ are the elementary symmetric functions in $t_1,\ldots,t_n$. [Cf. the Fundamental Theorem of Symmetric Polynomials 70.4] Let $D = \prod_{i<j}(t_i - t_j)$ in $\mathbb{Z}[t_1,\ldots,t_n]$. If $\sigma$ is the transposition $(ij)$ in $S_n$, then $\sigma D = -D$. It follows if $\sigma$ lies in $S_n$, then

$$\sigma(D) = \begin{cases} -D & \text{if } \sigma \notin A_n \\ D & \text{if } \sigma \in A_n. \end{cases}$$

Let $F$ be a field of characteristic different from two, $f$ a polynomial in $F[t]$, and $K/F$ a splitting field of $f$. Suppose that $\deg f = n$ and $f$ has no multiple roots in $K$, say the roots are $\alpha_1,\ldots,\alpha_n$. Hence $K/F$ is Galois and the Galois group $G(K/F)$ permutes $\alpha_1,\ldots,\alpha_n$, so defines a monomorphism $G(K/F) \to S_n$. If $\sigma$ lies in $G(K/F)$, we have a commutative diagram:

$$\begin{array}{c}
\frac{F[t_1,\ldots,t_n]}{e_{\alpha_1,\ldots,\alpha_n}} \xrightarrow{\bar{\sigma}} \frac{F[t_1,\ldots,t_n]}{e_{\alpha_1,\ldots,\alpha_n}} \\
\sigma \downarrow \quad \quad \downarrow \\
K \quad \quad K,
\end{array}$$
where \( \hat{\sigma} \) is the automorphism of \( F[t_1, \ldots, t_n] \) induced by \( t_i \mapsto t_{\sigma(i)} \) for \( i = 1, \ldots, n \), and \( e_{\alpha_1, \ldots, \alpha_n} \) is the obvious evaluation map. It follows that
\[
d := e_{\alpha_1, \ldots, \alpha_n}(D) = \prod_{i<j}(\alpha_i - \alpha_j) \neq 0
\]
satisfies
\[
\sigma(d) = \begin{cases} -d & \text{if } \sigma \notin A_n \\ d & \text{if } \sigma \in A_n. \end{cases}
\]

Therefore, \( G(K/F(d)) = G(K/F) \cap A_n \). Let
\[
\Delta := d^2 = \prod_{i<j}(\alpha_i - \alpha_j)^2.
\]
Then \( \sigma(\Delta) = \Delta \) for every \( \sigma \in G(K/F) \), hence \( d^2 = \Delta \) lies in \( K^{G(K/F)} = F \) and is called the discriminant of \( f \). As
\[
[G(K/F) : G(K/F(d))] \leq [S_n : A_n] = 2,
\]
either \( F = F(d) \) or \( F(d)/F \) is of degree two, as \( d \in F \) if and only if \( \Delta \in F^2 \) if and only if \( G(K/F) \) is a subgroup of \( A_n \) (where we view \( G(K/F) \to S_n \) as an inclusion).

**Computation 63.1.** In the setup above, let \( s_i = s_i(\alpha_1, \ldots, \alpha_n) \). As \( S_n \) fixes each \( s_j(t_1, \ldots, t_n) \), we have each \( s_i \) lies in \( K^{G(K/F)} = F \). By linear algebra, the Vandermonde determinant
\[
\det \begin{pmatrix} 1 & \cdots & 1 \\ \alpha_1 & \cdots & \alpha_n \\ \vdots & \vdots & \vdots \\ \alpha_1^{n-1} & \cdots & \alpha_n^{n-1} \end{pmatrix} = \prod_{i>j}(\alpha_i - \alpha_j) = d,
\]
so if
\[
A = \begin{pmatrix} 1 & \cdots & 1 \\ \alpha_1 & \cdots & \alpha_n \\ \vdots & \vdots & \vdots \\ \alpha_1^{n-1} & \cdots & \alpha_n^{n-1} \end{pmatrix},
\]

is the Vandermonde matrix, we have \( \det AA^t = d^2 = \Delta \).

**Check 63.4.** Let \( e_i := \sum_{j=1}^n \alpha_j^i \) for \( i = 1, \ldots, 2n - 2 \). Then
\[
AA^t = \begin{pmatrix} n & e_1 & e_2 & \cdots & e_{n-1} \\ e_1 & e_2 & e_3 & \cdots & e_n \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ e_{n-1} & e_n & e_{n+1} & \cdots & e_{2n-2} \end{pmatrix}.
\]
As \( \sum_{j=1}^{n} t_j^i \) lies in \( F[t_1, \ldots, t_n]^{S_n} \), for every positive integer \( i \), each \( \sum_{j=1}^{n} t_j^i \) is a polynomial in the \( s_i(t_1, \ldots, t_n) \)'s (in fact, with integer coefficients), hence each \( e_i \) lies in \( F[s_1, \ldots, s_n] \). We can, therefore, compute formulas for \( \Delta \) for specific \( n \).

**Example 63.5.** In the computation and the notation there, we let \( n = 2 \). Then

\[
f = (t - \alpha_1)(t - \alpha_2) = t^2 - (\alpha_1 + \alpha_2)t + \alpha_1\alpha_2 = t^2 - s_1t + s_2.
\]

So \( s_1 = \alpha_1 + \alpha_2 \) and \( s_2 = \alpha_1\alpha_2 \), hence \( e_1 = s_1 = \alpha_1 + \alpha_2 \) and \( e_2 = \alpha_1^2 + \alpha_2^2 = (\alpha_1 + \alpha_2)^2 - 2\alpha_1\alpha_2 = s_1^2 - 2s_2 \). It follows that

\[
\Delta = \det \begin{pmatrix} 2 & e_1 \\ e_2 & e_2 \end{pmatrix} = 2e_2 - e_1^2 = 2(s_1^2 - 2s_2) - s_1^2 = s_2^2 - 4s_2.
\]

In general, if \( f \) is a polynomial of degree \( n \) and splits over \( K \) with the roots \( \alpha_1, \ldots, \alpha_n \) (possibly not distinct), we define the discriminant of \( f \) to be \( \Delta = \prod_{i < j}(\alpha_i - \alpha_j)^2 \).

If one knows a root of an irreducible polynomial, then we can express the discriminant of it in another way.

**Proposition 63.6.** Let \( f \) be an irreducible and separable polynomial in \( F[t] \) of degree \( n \), \( L/F \) a splitting field of \( f \), \( \alpha = \alpha_1, \alpha_2, \ldots, \alpha_n \) the (distinct) roots of \( f \) in \( L \), \( K = F(\alpha) \), and \( T_\alpha : L \rightarrow L \) the \( F \)-linear isomorphism given by \( x \mapsto \alpha x \). Then the following are true:

1. \( \alpha_1, \ldots, \alpha_n \) are the eigenvalues of \( T_\alpha \). In particular, \( T_\alpha \) is diagonalizable.
2. \( \alpha_1^k, \ldots, \alpha_n^k \) are the eigenvalues of \( T_\alpha^k \) for all positive integers \( k \) and \( T_\alpha^k \) is diagonalizable.
3. If \( f = t^n + a_{n-1}t^{n-1} + \cdots + a_0 \) in \( F[t] \), then \( f_{T_\alpha} = f_{[L:K]} \), where \( f_{T_\alpha} \) is the characteristic polynomial of \( T_\alpha \). In particular,
   a. \( \text{Tr}_{K/F}(\alpha) = -a_{n-1} \).
   b. \( N_{K/F}(\alpha) = (-1)^n a_0 \).
4. Let \( A \) be the Vandermonde matrix 63.3. Then the \((ij)\)th term of \( AA^t \) satisfies
   \[
   (AA^t)_{ij} = \text{Tr}_{K/F}(\alpha^{i-1} \alpha^{j-1}).
   \]
   In particular, \( \det \left( \text{Tr}_{K/F}(\alpha^{i-1} \alpha^{j-1}) \right) = \Delta(\alpha) \), hence nonzero.
5. We have
   \[
   \Delta(\alpha) = (-1)^{\frac{n(n-1)}{2}} N_{K/F}(f'(\alpha)).
   \]

**Proof.** (1) and (2) are left as new exercises and (3) is Exercise 59.26(3).
(4): Using (2), we have $\text{Tr}_{K/F}(\alpha^k) = \alpha_1^k + \cdots + \alpha_n^k$ for all $k \geq 1$, hence

$$(AA^t)_{ij} = \sum_{l=1}^{n} \alpha_i^{l-1} \alpha_j^{l-1} = \sum_{l=1}^{n} \alpha_i^{l} \alpha_j^{l-2} = \text{Tr}_{K/F}(\alpha_i^{-1} \alpha_j^{l-1}).$$

(5): We have $\det A = \prod_{i>j}(\alpha_i - \alpha_j)$ and there exist $n(n-1)/2$ pairs of integers $(i, j)$ with $1 \leq j < i \leq n$ in the product $\Delta(\alpha) = \prod_{i>j}(\alpha_i - \alpha_j)^2$. As $(\alpha_i - \alpha_j)^2 = -(\alpha_i - \alpha_j)(\alpha_j - \alpha_i)$, it follows that

$$\Delta(\alpha) = \left(-1\right)^{\frac{n(n-1)}{2}} \prod_{i=1}^{n} \prod_{j=1, j\neq i}^{n} (\alpha_i - \alpha_j).$$

Since $f = \prod_{i=1}^{n}(t - \alpha_i)$ in $L[t]$, we have

$$f' = \sum_{i=1}^{n} \prod_{j=1, j\neq i}^{n} (t - \alpha_j), \quad \text{so} \quad f'(\alpha_i) = \prod_{j=1, j\neq i}^{n} (\alpha_i - \alpha_j)$$

and

$$\Delta(\alpha) = \left(-1\right)^{\frac{n(n-1)}{2}} \prod_{i=1}^{n} f'(\alpha_i).$$

Let $\sigma_i : K \rightarrow L$ be the $F$-homomorphisms satisfying $\alpha \mapsto \alpha_i$ for $i = 1, \ldots, n$. Then $\sigma_i(f'(\alpha)) = f'(\alpha_i)$, hence

$$\Delta(\alpha) = \left(-1\right)^{\frac{n(n-1)}{2}} N_{K/F}(f'(\alpha))$$

as claimed. \qed

We use the above in the following computation:

**Example 63.7.** In Proposition 63.6, set $F = \mathbb{Q}$ and $L = K = \mathbb{Q}(\zeta)$, where $\zeta$ a primitive $p$th root of unity in $\mathbb{C}$ with $p$ an odd prime. The $p$th cyclotomic polynomial $\Phi_p = t^p - 1 + \cdots + 1 = m_\mathbb{Q}(\zeta)$. So by Proposition 63.6(5), we have

$$\Delta(\zeta) = \left(-1\right)\frac{\left(p-1\right)(p-2)}{2} N_{K/F}(\Phi_p'(\zeta)) = \left(-1\right)^{\frac{p-1}{2}} N_{K/F}(\Phi_p'(\zeta)).$$

We compute the right hand side of this equation. Since $t^p - 1 = (t - 1)\Phi_p$, taking the derivative yields $pt^{p-1} = \Phi_p + (t - 1)\Phi'_p$. Hence

$$\Phi'_p(\zeta) = \frac{p\zeta^{p-1}}{\zeta - 1}.$$

By Proposition 63.6(5),

$$N_{\mathbb{Q}(\zeta)/\mathbb{Q}}(\Phi_p'(\zeta)) = \frac{N_{\mathbb{Q}(\zeta)/\mathbb{Q}}(p\zeta^{p-1})}{N_{\mathbb{Q}(\zeta)/\mathbb{Q}}(\zeta - 1)} = \frac{p^{p-1} \cdot 1}{N_{\mathbb{Q}(\zeta)/\mathbb{Q}}(\zeta - 1)}.$$
As \( \prod_{i=1}^{p-1} (t - \zeta^i) = t^{p-1} + \cdots + t + 1 \), we have

\[
\prod_{i=1}^{p-1} (\zeta^i - 1) = (-1)^{p-1} \Phi_p(1) = (-1)^{p-1} p = p.
\]

Consequently, \( N_{\mathbb{Q}(\zeta)/\mathbb{Q}} \left( \Phi'_p(\zeta) \right) = p^{p-2} \), hence

\[
\Delta(\zeta) = (-1)^{\frac{p-1}{2}} p^{p-2}.
\]

In particular, \( \mathbb{Q}(\sqrt{\Delta(\zeta)}) \) is the unique quadratic extension of \( \mathbb{Q} \) in \( \mathbb{Q}(\zeta) \) by Proposition 58.14.

Note, since

\[
\prod_{p-1 \geq i > j \geq 1} (\zeta^i - \zeta^j)^2 = \prod_{p-1 \geq i > j \geq 1} ((\zeta^i - 1) - (\zeta^j - 1))^2,
\]

we also have \( \Delta(\zeta) = \Delta(1 - \zeta) \).

The technique above can be generalized in order to compute \( \Delta(\zeta) \) with \( \zeta \) a primitive \( p^m \)th root of unity (and from it \( \zeta \) a primitive \( n \)th root of unity). To show this, one uses:

**Lemma 63.8.** Let \( \zeta \) be a primitive \( n \)th root of unity in \( \mathbb{C} \) and \( K = \mathbb{Q}(\zeta) \). Then \( \Delta(\zeta) \mid n^{\varphi(n)} \).

**Proof.** Let \( \Phi_n = m_{K/F}(\zeta) \), then \( t^n - 1 = \Phi_ng \) for some \( g \in \mathbb{Z}[t] \).

Taking derivatives, we have \( nt^{n-1} = \Phi'_ng + \Phi_ng \). Evaluating at \( t = \zeta \) and taking norms give

\[
n^{\varphi(n)} = N_{K/\mathbb{Q}}(n\zeta^{n-1}) = N_{K/\mathbb{Q}}(\Phi'_n) N_{K/\mathbb{Q}}(g) = (-1)^{\frac{n(n-1)}{2}} \Delta(\zeta) N_{K/\mathbb{Q}}(g)
\]

in \( \mathbb{Z} \). The result follows. \( \square \)

**Exercises 63.9.**

1. Show that equation 63.2 is valid.

2. Prove Check 63.4.

3. Let \( f = t^3 - a_1t^2 - a_2t - a_3 \in \mathbb{R}[t] \). Show
   
   (a) The discriminant \( \Delta = -4a_1^3a_3 + a_1^2a_2^2 - 18a_1a_2a_3 + 4a_2^3 - 27a_3^2 \).
   (b) \( f \) has multiple roots if and only if \( \Delta = 0 \).
   (c) \( f \) has three distinct real roots if and only if \( \Delta > 0 \).
   (d) \( f \) has one real root and two non-real roots if and only if \( \Delta < 0 \).

4. Let \( F \) be a field and \( f \) be an irreducible separable cubic in \( F[t] \). Show the Galois group of \( f \) is either \( A_3 \) or \( S_3 \) and if the characteristic of \( FE \) is not two, it is \( A_3 \) if and only if the discriminant of \( f \) is a square in \( F \).
5. Let \( t^3 + pt + q \) be irreducible over a finite field \( F \) of characteristic not 2 or 3. Show that \(-4p^3 - 27q^2\) is a square in \( F \).

6. Prove (1) and (2) of Proposition 63.6

7. Let \( \zeta \) be a primitive \( p^m \)-th root of unity in \( \mathbb{C} \) and \( K = \mathbb{Q}(\zeta) \). Show that \( \Delta(\zeta) = \pm p^{m-1}(mp-m-1) \) with the plus sign occurring if and only if either \( p \equiv 1 \mod 4 \) or \( a = 2^a \) with \( a > 2 \).

### 64. Purely Transcendental Extensions

Let \( K/F \) be a finitely generated extension of fields in which \( K \) is *purely transcendental* over \( F \), i.e., \( K \cong F(t_1, \ldots, t_n) \). If \( K/E/F \) is an intermediate field, the question arises whether \( E/F \) is a purely transcendental extension. In general, this has been shown to be false, but if \( n = 1 \) this is true and is the subject of this section.

**Lemma 64.1.** Let \( F \) be a field and \( x \) an element transcendental over \( F \). Suppose that \( u \) is an element in \( F(x) \) not lying in \( F \), say \( u = f/g \) with \( f \) and \( g \) nonzero non-constant relatively prime polynomials in \( F[t] \), and \( r = \max\{\deg f, \deg g\} \). Then \( x \) is algebraic over \( F(u) \) with \( m_{F(u)}(X) = f - ug \) in \( F(u)[t] \).

**Proof.** Let \( h = f - ug \) in \( F[u][t] \subset F(u)[t] \) with \( f \) and \( g \) as in the lemma. Suppose that \( f = \sum a_it^i \) and \( g = \sum b_it^i \). If \( b_j \) is nonzero then so is \( a_j - ub_j \), lest \( u \) lies in \( F \). Thus \( \deg h = \max\{\deg f, \deg g\} \) and \( x \) is algebraic over \( F(u) \). It follows that \( u \) is transcendental over \( F \) as \( x \) is. To finish, we must show that \( h \) is irreducible in \( F(u)[t] \). Suppose this is false, then \( h \) is also reducible in the UFD \( F[u, t] \) by a consequence of Gauss’ Lemma (Lemma 34.7). As \( h \) is primitive, we have \( h = h_1h_2 \) in \( F[u, t] \) with \( h_1 \) and \( h_2 \) lying in \( F[u, t] \setminus F \). As \( u \) occurs linearly in \( h \), we may assume that it occurs in \( h_1 \), say \( h_1 = uh_3 + h_4 \) with \( h_3, h_4 \) in \( F[t] \) and \( h_2 \) lies in \( F[t] \). So we have

\[
g = f - ug = uh_2h_3 - h_2h_4.
\]

As \( u \) is transcendental over \( F(t) \), \( \text{w} \) must have \( f = -h_2h_4 \) and \( g \) also. Therefore, \( h_2 | f \) and \( h_2 | g \) in \( F[t] \), a contradiction. \( \square \)

**Corollary 64.2.** Let \( F \) be a field, \( x \) a transcendental element over \( F \), and \( u \) an element of \( F(x) \). Then \( F(u) = F(x) \) if and only if there exist elements \( a, b, c, d \) in \( F \) with \( ad - bc \) nonzero and \( u = \frac{ax + b}{cx + d} \).

[Note the condition on \( a, b, c, d \) insures that \( u \) does not lie in \( F \).]
Proposition 64.3. Let $F$ be a field and $x$ a transcendental element over $F$. Then
\[
G(F(x)/F) = \{ f : F(x) \to F(x) \mid x \mapsto \frac{ax + b}{cx + d} \text{ an } F\text{-homomorphism with } a, b, c, d \text{ in } F \text{ satisfying } ad - bc \neq 0 \}
\]
\[
\cong \text{GL}_2(F)/\text{Z} \left( \text{GL}_2(F) \right).
\]
The group $\text{GL}_2(F)/\text{Z} \left( \text{GL}_2(F) \right)$ is called the (second) projective linear group and denoted by $\text{PGL}_2(F)$. If $F = \mathbb{C}$, then $\mathbb{C} = \mathbb{C}^2$, so $\text{PGL}_2(\mathbb{C}) = \text{PSL}_2(\mathbb{C}) := \text{SL}_2(F)/\text{Z} \left( \text{SL}_2(F) \right)$. (If $A \in \text{GL}_2(\mathbb{C})$ multiply it by $\begin{pmatrix} d & 0 \\ 0 & d \end{pmatrix}$ where $d = (\sqrt{\det A})^{-1}$.)

Theorem 64.4. (Lüroth’s Theorem) Let $F$ be a field and $x$ a transcendental element over $F$. If $F(x)/E/F$ is an intermediate field, then there exists an element $y$ in $F(x)$ such that $E = F(y)$.

Proof. We may assume that $F < E \subset F(x)$. Let $z$ be an element in $E$ not in $F$. Then $F(x)/F(z)$ is algebraic by the lemma, hence $F(x)/E$ is algebraic. Let $f = a m_E(x)$ with $a \in F[x]$ and $f \in F[x][t]$ primitive. Say
\[
f = a_n(x) t^n + \cdots + a_0(x) \text{ with } a_i \in F[x] \text{ for all } i
\]
where $n$ is the degree of $f$ in $t$, so $a_n(x)$ is nonzero and $n = \deg m_E(x) = |F(x) : E|$. We also have all $a_i(x)/a_n(x)$ lie in $E$ but as $x$ is transcendental over $F$, there exists an $i$ such that $u = a_i(x)/a_n(x)$ does not lie in $F$. We can also write $u = g(x)/h(x)$ with $g$ and $h$ nonzero relatively prime polynomials in $F[t]$. If $r = \max \{ \deg g, \deg h \}$ then $r = |F(x) : F(u)|$ by the lemma, so
\[
n = |F(x) : E| \leq |F(x) : F(u)| = r.
\]
Let $d = g(x)h(t) - h(x)g(t)$ in $F[x,t]$. As $g$ and $h$ are relatively prime in the UFD $F[x,t]$, the polynomial $d$ is nonzero. In addition, the polynomial
\[
uh(t) - g(t) = h(x)^{-1}d \text{ in } E[t] \text{ has } x \text{ as a root,}
\]
so $m_E(x) \mid h(x)^{-1}d$ in $E[t]$, hence $f = a(x)m_E(x) \mid h(x)^{-1}d$ in $F(x)[t]$. Since $f$ is primitive in $F[x,t]$ and $d$ lies in $F[x,t]$, we must have $f \mid d$ in $F[x,t]$, as a consequence of Gauss’ Lemma (e.g., by Exercise 3ii(34.12)). Write $d = \alpha f$ with $\alpha \in F[x,t]$ and let $\deg_x$ denote the degree in $x$. We have $\deg_x d \leq r$ and $\deg_x f \geq \max \{ \deg a_j(x) \mid a_j(x) \neq 0 \}$. By choice, $u = g(x)/h(x) = a_i(x)/a_n(x)$, so $\deg_x f \geq \max \{ \deg g, \deg h \} =$
r. Consequently, \( d = \alpha f \) in \( F[x,t] \) and \( f \mid d \) in \( F[x,t] \) implies that \( \deg_x d = r = \deg_x f \) and \( \deg_x \alpha = 0 \). In particular, \( \alpha \) lies in \( F[t] \), so is primitive when viewed as a polynomial over \( F[x] \), i.e., an element in \( F[x][t] \). By Gauss’ Lemma, \( d = f \alpha \) is also primitive in \( F[x][t] \). Since \( \alpha \) lies in \( F[t] \) and \( \alpha \mid d \) in \( F[x,t] \), we must have \( \alpha \) lies in \( F[t]^\times = F^\times \), hence
\[
n = \deg_t f = \deg_t d = \deg_x d = \deg_x f = r,\]
so \( E = F(u) \). \( \square \)

It is known if \( F \) is an algebraically closed field of characteristic zero, e.g., \( \mathbb{C} \), that any subfield of \( F[t_1,t_2] \) is purely transcendental. [A similar result holds if the characteristic is positive with a mild additional assumption.] However, the result is false for \( F[t_1,t_2,t_3] \).

Suppose that \( F \) is a field and \( x \) is transcendental over \( F \). Let \( H \subset G(F(x)/F) \cong \text{PGL}_2(F) \) be a finite subgroup. Then \( F(x)/F(x)^H \) is a finite Galois extension. By Lüroth’s Theorem, there exists an element \( u \) in \( F(x) \) such that \( F(u) = F(x)^H \). Suppose that \( F = \mathbb{C} \), then using a bit of complex analysis (stereographic projections, etc.), one can show that if \( H \subset \text{PGL}_2(\mathbb{C}) \) is a finite subgroup, then \( H \) is isomorphic to a finite rotation group in \( \mathbb{R}^3 \). In Section 20 we classified such rotation groups in \( \mathbb{R}^3 \). Recall that they are the following:

1. cyclic: \( \cong \mu_n \), all \( n \) (planar).
2. dihedral: \( \cong D_n \), all \( n \) (planar).
3. tetrahedral: \( \cong A_4 \) (the rotations of a regular tetrahedron).
4. octahedral: \( \cong S_4 \) (the rotations of a cube or a regular octahedron).
5. icosahedral: \( \cong A_5 \) (the rotations of a regular dodecahedron or a regular icosahedron).

**Exercise 64.5.** Prove that the rotations of a regular tetrahedron is isomorphic to \( A_4 \), the rotations of a cube or octahedron is isomorphic to \( S_4 \), and the rotations of a dodecahedron or icosahedron is isomorphic to \( A_5 \) using Exercise 24.24(7).

### 65. Finite Fields

A major property of a finite field \( F \) is that the multiplicative group of \( F \) is cyclic. If you have done the appropriate exercises, you see that we have all of the following:

**Theorem 65.1.** Let \( F \) be a finite field with \( |F| = q \). Then all the following are true:
(1) $F^\times$ is a cyclic group.
(2) The characteristic of $F$ is $p$ for some prime $p$ and $q = p^n$.
(3) If $K/F$ is a finite extension of fields of degree $m$, then $|K| = q^m$.
(4) If $p$ is a prime and $m$ a positive integer, then there exists a field with $p^m$ elements.
(5) If $|F| = p^n$, with $p$ a prime, then $F$ is a splitting field of $t^{p^n} - t$ over $\mathbb{Z}/p\mathbb{Z}$.
(6) Any two finite fields with the same number of elements are isomorphic. In particular, in a fixed algebraic closure of $\mathbb{Z}/p\mathbb{Z}$, with $p$ a prime, there exists a unique field of order $p^m$ for each positive integer $m$.
(7) If $|F| = q = p^n$, then $\varphi_F : F \to F$ by $x \mapsto x^q$ is an $F$-automorphism, the Frobenius automorphism, and $G(F/\mathbb{Z}/p\mathbb{Z}) = \langle \varphi_F \rangle$.
(8) If $|F| = p^n$ and $K$ a field with $|K| = p^m$ in an algebraic closure of $F$, then $K/F$ if and only if $n \mid m$.
(9) If $K/F$ is a finite extension of degree $m$ and $\varphi_K : K \to K$ the Frobenius automorphism $x \mapsto x^{q^m}$, then $K/F$ is a cyclic extension with $G(K/F) = \langle \varphi_K \rangle$.
(10) Any finite field is perfect.
(11) There exists an $x$ in $F$ such that $F = (\mathbb{Z}/p\mathbb{Z})(x)$.

Given a finite field $F$ with $q = p^n$ elements, one would like to find irreducible polynomials in $F[t]$. If $f$ is an irreducible polynomial of degree $d$ in $F[t]$, then by the theorem, $f \mid t^{p^n} - t$ if and only if $d \mid n$. In particular, we must have

$$t^{p^n} - t = \prod_{d \mid n} f_{i,d}$$

where the $f_{i,d}$ in $F[t]$ run over all the monic irreducible polynomials of degree $d$. So we have $p^n = \sum_{d \mid n} dN_{p,d}$, where $N_{p,d}$ is the number of monic irreducible polynomials of degree $d$ in $F[t]$, by comparing degrees. By the general form of Mönbius Inversion (Exercise 58.21(6)), we conclude that

**Proposition 65.2.** The number $N_{p,n}$ of monic irreducible polynomials of degree $n$ in $(\mathbb{Z}/p\mathbb{Z})[t]$, with $p$ a prime, satisfies

$$N_{p,n} = \frac{1}{n} \sum_{d \mid n} \mu\left(\frac{n}{d}\right)p^d.$$

Using the general form of the Mönbius Inversion given by Exercise 58.21(6), we can determine irreducible polynomials, just as in the case
of cyclotomic extensions of \( \mathbb{Q} \). Let \( \zeta \) be a primitive root of unity over \( \mathbb{Z}/p\mathbb{Z} \), i.e., a generator for \( F^\times \) if \( F \) is a field with \( p^n \) elements satisfying \( p \nmid n \). Let \( \Phi_n \) denote the product of \( t - \zeta \) where \( \zeta \) is a primitive \( n \)-th root of unity in an algebraic closure of \( \mathbb{Z}/p\mathbb{Z} \), so \( \deg \Phi_n = \phi(n) \). Unlike the case of \( \mathbb{Q} \), the polynomial \( \Phi_n \) is not irreducible. We have \( t^n - 1 = \prod_{d \mid n} \Phi_d \). Arguing as the case over \( \mathbb{Q} \), we conclude that

\[
\Phi_n = \prod_{d \mid n} \Phi_d = \left( \prod_{d \mid n} (t^d - t) \right)^{\mu(d)},
\]

in \( (\mathbb{Z}/p\mathbb{Z})[t] \) by Möbius Inversion.

**Example 65.3.** Let \( p \) be a prime not dividing 12. Then

\[
\Phi_{12} = \prod_{d \mid 12} (t^{\frac{12}{d}} - 1)^{\mu(d)} = (t^6 - t^4 + t^2 - 1)(t^3 - 1)^2(t^2 - 1)^3(t - 1)^6.
\]

**Lemma 65.4.** Let \( F \) be a finite field with \( q \) elements and \( i \) a positive integer. Set

\[
S(t) := \sum_{x \in F} x^i \text{ in } F,
\]

then

\[
S(t) = \begin{cases} 
-1 & \text{if } q - 1 \mid i \\
0 & \text{otherwise.}
\end{cases}
\]

For convenience, as \( q = 0 \) in \( F \), we shall set \( S(t^0) = 0 \) also, i.e., let \( 0^0 = 1 \).

**Proof.** Suppose that \( q - 1 \mid i \). Then \( x^i = 1 \) for all \( x \) in \( F^\times \) and

\[
S(t) = \sum_{x \in F^\times} x^i = |F^{q-1}| 1_F = -1 \text{ in } F.
\]

So we may assume that \( q - 1 \nmid i \). Then there exists an element \( y \) in \( F^\times \) such that \( y^i \neq 1 \) as \( F^\times \) is cyclic and the gcd of \( |F^\times| \) and \( i \) cannot be \( i \). Since \( yF^\times = F^\times \), we have

\[
S(t) = \sum_{x \in F^\times} x^i = \sum_{x \in F^\times} (yx)^i = y^i \sum_{x \in F^\times} x^i,
\]

so \( (1 - y^i) \sum_{x \in F^\times} x^i = 0 \). It follows that \( S(t^0) = 0 \).
max\{i_1 + \cdots + i_n \mid a_{i_1,\ldots,i_n} \text{ nonzero}\}. We are interested in solutions of \( f = 0 \) in \( R^n \) for such an \( f \) when \( R \) is a finite field.

**Theorem 65.5.** (Chevalley-Warning Theorem) Let \( F \) be a finite field of characteristic \( p \) and \( f \) a non-constant polynomial in \( F[t_1,\ldots,t_n] \) having total degree \( \deg f < n \). Set \( V = \{ x = (x_1,\ldots,x_n) \in F^n \mid f(x) = 0 \text{ in } F \} \). Then \( |V| \equiv 0 \mod p \).

**Proof.** Suppose that \( F \) has \( q \) elements. Set \( P = 1 - f^{q-1} \) in \( F[t_1,\ldots,t_n] \). As \( |F^\times| = q - 1 \), if \( x \) lies in \( V \), then \( f(x) = 0 \) in \( F \), so \( P(x) = 1 \) in \( F \); and if \( x \) does not lie in \( V \), then \( P(x) = 0 \) in \( F \), i.e.,

\[
P(x) = \begin{cases} 0 & \text{if } x \notin V \\ 1 & \text{if } x \in V. \end{cases}
\]

(So \( P \) is the characteristic function for \( V \).) For any \( g \) in \( F[t_1,\ldots,t_n] \), let

\[
S(\hat{g}) := \sum_{F^n} g(x) \text{ in } F,
\]

so

\[
S(\hat{P}) = \sum_{F^n} P(x) = \sum_{V} 1 = |V|1_F \text{ in } F.
\]

In particular, if we show \( S(\hat{P}) = 0 \) in \( F \), then \( |V| \equiv 0 \mod p \) as needed.

We know that for all \((i_1,\ldots,i_n)\),

\[
\sum_{F^n} x_1^{i_1} \cdots x_n^{i_n} = (\sum_{F} x_1^{i_1}) \cdots (\sum_{F} x_n^{i_n})
\]

and \( \deg P < n(q-1) \). If \( P = \sum_{i_1,\ldots,i_n} a_{i_1,\ldots,i_n} t_1^{i_1} \cdots t_n^{i_n} \), we have \( \sum_{j=1}^n i_j < n(q-1) \) if \( a_{i_1,\ldots,i_n} \neq 0 \). In particular, if \( a_{i_1,\ldots,i_n} \neq 0 \) is a coefficient of \( P \), then, there exists a \( j, 1 \leq j \leq n \), with \( i_j < q - 1 \). Therefore, we have

\[
S(\hat{P}) = \sum_{F^n} P(x) = \sum_{F^n} \sum_{i_1,\ldots,i_n} a_{i_1,\ldots,i_n} x_1^{i_1} \cdots x_n^{i_n}
\]

\[
= \sum_{i_1,\ldots,i_n} a_{i_1,\ldots,i_n} \left( \sum_{F} x_1^{i_1} \right) \cdots \left( \sum_{F} x_n^{i_n} \right)
\]

\[
= \sum_{i_1,\ldots,i_n} a_{i_1,\ldots,i_n} \left( \sum_{F} x_1^{i_1} \right) \cdots \left( \sum_{F} x_n^{i_n} \right) = 0,
\]

as \( \sum_{F} y^j = S(\hat{t}^j) = 0 \) for \( j < q - 1 \) by the lemma. \( \square \)
Of course, one is interested in solutions of such a non-constant polynomial \( f \) in \( F[t_1, \ldots, t_n] \) if every variable \( t_i \) occurs in \( f \) nontrivially, since the result is always true otherwise. A polynomial \( f \) in \( F[t_1, \ldots, t_n] \) is called homogeneous of degree \( d \) if every nonzero monomial in \( f \) has the same total degree \( d \).

**Corollary 65.6.** Let \( F \) be a finite field and \( f \) a non-constant polynomial in \( F[t_1, \ldots, t_n] \) having total degree \( \deg f < n \). Suppose that \( f(0) = 0 \). Then there exists a point \( x \) in \( F^n \) besides the origin such that \( f(x) = 0 \). In particular, if \( f \) is homogeneous of degree \( d \) satisfying \( n > d \), then \( f \) has a nontrivial solution in \( F^n \).

A homogeneous polynomial of degree two is called a quadratic form. The corollary says that any quadratic form in at least three variables over a finite field has a nontrivial solution. A much deeper theorem (due to Weil) says that if \( f \) is any homogeneous polynomial in \( \mathbb{Z}[t_1, \ldots, t_n] \) that remains irreducible in \( \mathbb{C}[t_1, \ldots, t_n] \), then \( f \equiv 0 \mod p \) has a nontrivial zero for all primes \( p \gg 0 \).

A theorem of Jacobson says that if \( R \) is a ring (even a rng) in which for each element \( x \) there exists an integer \( n = n(x) > 1 \) satisfying \( x^n = x \), then \( R \) is commutative. For a proof of this more general result, see Section 66. The following result is the starting point.

**Theorem 65.7.** (Wedderburn) Every finite division ring is a field.

**Proof.** (Witt) Let \( D \) be a division ring. Define the center \( Z(D) \) of \( D \) to be \( Z(D) := \{x \in D \mid xy = yx \text{ for all } y \text{ in } D\} \). So \( Z(D^\times) \) is the center of the group \( D^\times \). Clearly, \( Z(D) \) is a field. Now suppose that \( D \) is finite. Let \( F = Z(D) \). Then \( F \) is also finite, so char \( F = p \) for some prime \( p \) and \( |F| = q = p^m \), for some \( n \). We know that \( D \) is an \( F \)-vector space, so \( |D| = q^m \) where \( m = \dim_F D \). Assume that \( D \) is not commutative, i.e., \( m > 1 \). If \( x \in D \) does not lie in \( F \), then the conjugacy class of \( x \) in \( D^\times \) is \( \{y \in D^\times \mid y = axa^{-1} \text{ for some } a \in D^\times \} \) and the centralizer of \( x \) in \( D^\times \) is \( Z_{D^\times}(x) = \{y \in D^\times \mid xy = yx\} \). Set \( Z_D(x) = Z_{D^\times}(x) \cup \{0\} \). Then \( F \subset Z_D(x) \subset D \). Clearly, \( Z_D(x) \) is also a division ring and a finite dimensional vector space over \( F \). Set \( \delta(x) := \dim_F Z_D(x) \). We have \( |Z_D(x)| = q^{\delta(x)} \). The notation and basic theorems about vector spaces over fields hold over division rings (i.e., as (left) modules over division rings) with the same proofs. [The one exception is the Representation Theorem of linear operators where matrix multiplication of matrices \( A \) and \( B \) must be modified so that the \( ij \)the entry of \( AB \) is \( \sum B_{kj}A_{ik} \) unless linear operators and scalars are written on opposite sides.] In particular, \( D \) is a (left) \( Z_D(x) \)-vector...
space. It is easy to see that \( q^m = |D| = |Z_D(x)|^r = q^{\delta(x)r} \), some \( r \). In particular, \( m = \delta(x)r \), so \( \delta(x) \mid m \) (cf. with the analogue of a tower of fields). By assumption \( \delta(x) < m \) as \( x \) does not lie in \( F \). Let \( C^* \) be a system of representatives of the conjugacy classes of \( D^\times \) not in the center of \( D^\times = F^\times \). By the Class Equation 21.3, we have

\[
q^m - 1 = |F^\times| + \sum_{C^*} [D^\times : Z_{D^\times}(x)] = (q - 1) + \sum_{C^*} \frac{q^m - 1}{q^{\delta(x)} - 1}.
\]

Let \( \Phi_m \) be the \( m \)th cyclotomic polynomial \( \prod_{\zeta \in \mu_m \text{ primitive}} \zeta \) in \( \mathbb{Z}[t] \). Then we know that \( \Phi_m \mid \frac{t^m - 1}{t^\delta - 1} \) in \( \mathbb{Z}[t] \) if \( \delta \mid m \) and \( \delta \neq m \), as \( t^m - 1 = \prod_{d \mid m} \Phi_d \). Therefore,

\[
\Phi_m(q) \mid q^m - 1 \quad \text{and} \quad \Phi_m(q) \mid \sum_{C^*} \frac{q^m - 1}{q^{\delta(x)} - 1}
\]

in \( \mathbb{Z} \). Consequently, \( \Phi_m(q) \mid q - 1 \) in \( \mathbb{Z} \).

**Claim.** If \( \zeta \) is a primitive \( m \)th root of unity, then \( |q - \zeta| > |q - 1| \):

This is clear, if you draw a picture or if \( \zeta = e^{2\pi \sqrt{-1}r/m} \) with \( r \) and \( m \) relatively prime, then

\[
|q - \zeta|^2 = (q - e^{2\pi \sqrt{-1}r/m})(q - e^{-2\pi \sqrt{-1}r/m}) = q^2 + 1 - 2q \cos \frac{2\pi r}{m} > q^2 + 1 - 2q = (q - 1)^2.
\]

As \( -1 < \cos \frac{2\pi r}{m} < 1 \), the claim is established. The claim implies that

\[
|\Phi_m(q)| = \prod_{\zeta \in \mu_n \text{ primitive}} |q - \zeta| > |q - 1|.
\]

Hence \( \Phi_m(q) \mid q - 1 \) in \( \mathbb{Z} \) is impossible. This contradiction shows that \( D = F \).

**Exercises 65.8.**

1. Prove Theorem 65.1.
2. Show that the \( n \)th cyclotomic polynomial \( \Phi_n \in \mathbb{Z}/p\mathbb{Z}[t] \) with \( p \) a prime not dividing \( n \) is a product of \( \varphi(n)/d \) of irreducible polynomials of the same degree \( d \).
3. Show that over over finite field \( F \), every element of \( F \) is a sum of two squares.
4. Let \( K/F \) be a finite extension of finite fields. Prove that the norm map \( N_{K/F} : K \to F \) is surjective.
5. Let \( F \) be a finite field and \( f \) a homogeneous form of degree 2 (i.e., a quadratic form) in two variables with both variables occurring non-trivially. Then every nonzero element \( x \) in \( F \) is a value of \( f \), i.e., there exists a \( y \) in \( F \) satisfying \( f(y) = x \).

66. Addendum: Jacobson’s Theorem

In this section, we generalize Wedderburn’s Theorem that finite division rings are fields.

Let \( R \) be any ring. Then its center, \( Z(R) := \{ x \in R \mid xy = yx \text{ for all } y \text{ in } R \} \), is a commutative subring. If \( D \) is a division ring, then \( F[x] \) is a commutative subring of \( D \) and its quotient field \( F(x) \) is a field and also lies in \( D \). Moreover, if \( x \) is algebraic over \( F \), i.e., \( F[x] \) is a finite \( F \)-vector space, then \( F(x) = F[y] \) and \( F(x)/F \) is finite. We have looked at two cases for \( D \), the case that \( D \) is finite and the case that \( D \) is the Hamiltonian quaternions (where it is a finite dimensional \( \mathbb{R} \)-vector space). In the first case, \( D \) was commutative, but in the second case this might not be true. When \( D \) is a finite division ring, we have \( x^{|D|} = 1 \) for all \( x \) in \( D \). Jacobson’s Theorem says that a ring \( R \) that satisfies \( x^n = x \) for each \( x \) in \( R \) for some integer \( n = n(x) > 1 \) is commutative. We prove this in this section. We also show that any such ring is a subring of a product of fields, each algebraic over a finite field.

If \( D \) is a division ring, then two nonzero elements \( x \) and \( y \) in \( D \) commute if and only if \( yxy^{-1} = x \). We want to exploit this in the special case of Jacobson’s Theorem when \( R \) is a division ring (and reduce to the finite division ring case that we have answered).

**Proposition 66.1.** Let \( D \) be a division ring that satisfies \( x^n = x \) for each \( x \) in \( D \) and some integer \( n > 1 \) depending on \( x \). Then \( D \) is a field.

**Proof.** We know that the center of \( D \) contains a prime field, so \( n1_D \) lies in the center for all integers \( n \). Let \( x \) be a nonzero element of \( D \). Then there exist positive integers \( n \) and \( n' \) such that \( x^n = x \) and \( (2x)^{n'} = 2x \). It follows with \( N = (n-1)(n'-1)+1 \) that \( (2^N-2)x = 0 \) in \( D \), so \( 2^N - 2 = 0 \). We conclude that the center of \( D \) has characteristic \( p \) for some prime \( p \). Let \( F \) be the prime field in the center of \( D \), so \( F \cong \mathbb{Z}/p\mathbb{Z} \). Suppose that \( D \) is not commutative. Let \( x \) in \( D \) not lie in the center. As \( x \) is algebraic over \( F \), since it is a root of \( t^n - t \) in \( F[t] \) for some positive integer \( n \), the field extension \( F(x)/F \) is finite by the discussion prior to the proposition. Therefore, \( F(x) \) is a finite field. Indeed the same argument shows that \( F(y) \) is a finite field for all \( y \) in
Let $|F(x)| = p^m$. We also know that any element $z$ in $F(x)$ satisfies $z^{p^m} = z$.

**Claim:** There exists an element $y$ in $D$ and an integer $k$ such that $x^k \neq x$ and $yxy^{-1} = x^k$.

Let $T_x : D \to D$ be the $F$-linear transformation defined by $z \mapsto zx - xz$. As $x$ is not in the center of $D$, this map is not the trivial map, i.e., $T_x(z) \neq 0$ for some $z$. Since

$$T^2_x(z) = T_x(zx - xz) = zx^2 - 2xzx + x^2z,$$

inductively, we see that

$$T^l_x(z) = \sum_{i=0}^{l} (-1)^i \binom{l}{i} x^i zx^{l-i}$$

for all positive integers $l$ and all $z$ in $D$.

Since $\text{char } F = p$, the Children’s Binomial Theorem implies that $T^p_x(z) = zx^p - x^pz$, hence

$$T^p_x(z) = zx^p - x^p z = z(x^p - z)$$

for all $z$ in $D$, i.e., $T^p_x = T_x$. For each $w$ in $D$, left multiplication by $w$, $\lambda_w : D \to D$, i.e., $z \mapsto wz$ is also an $F$-linear transformation. We have for each $w$ in $F(x)$ (so $wx = xw$)

$$T_x \lambda_w(z) - T_x(wz) = wz - x(wz) = w(zx - xz) = \lambda_w T_x(z),$$

i.e., $T_x \lambda_w = \lambda_w T_x$ for all $w$ in $F(x)$ and $z$ in $D$. As $\lambda_w \lambda_{w'} = \lambda_{w'w}$ if $w, w'$ lie in $F(x)$, we have

$$0 = T^p_x - T_x = \prod_{F(x)} (T_x - \lambda_w).$$

By assumption, there exists an element $y$ in $D$ satisfying $xy - yx \neq 0$. As $D$ is a division ring, $D^\times$ satisfies the cancellation law and $y$ is invertible in $D$. It follows that there exists a nonzero element $w$ in $F(x)$ such that $0 = (T_x - \lambda_w)(y) = yx - xy - wy$, i.e., $yxy^{-1} = x + w$ in $F(x)$ by (*). Let $s = p^m - 1$, the order of $F(x)^\times = \langle x \rangle$. As $1, x, \ldots, x^{s-1}$ are all the roots of $t^s - 1$ in $F(x)[t]$, we must have $yxy^{-1} = x^k$ for some $k > 1$. This proves the claim.

Let $y$ in $D$ be as in the claim. We know by the argument above that $F(y)$ is also a finite field with say $p^n$ elements. So we have $y^{p^n} = y$ and $xy \neq yx$. Set

$$W := \{ z \in D \mid z = \sum_{i=0}^{p^m} \sum_{j=0}^{p^n} a_{ij}x^i y^j \text{ with } a_{ij} \in F \}.$$
We have $W$ a finite set closed under addition. By the claim, $W$ is closed under multiplication, so $W$ is a finite subring of $D$. If $z$ lies in $W$, $F(z) = F[z] \subseteq W$, so $W$ is a finite division ring. Consequently, $W$ is a field by Wedderburn’s Theorem, contradicting $x$ and $y$ do not commute.

Let $F$ be a field and $D$ a division ring containing $F$. If $x \in D$, then $F(x)$ is a field in $D$. We say that $D$ is algebraic over $F$ if every element in $D$ is algebraic over $F$.

Corollary 66.2. Let $D$ be a division algebra that is algebraic over a finite field. Then $D$ is commutative. In particular, $D/F$ is an algebraic extension of fields.

Proof. If $F$ is the finite field in $D$ and of characteristic $p$, then $F(x)$ is a finite field for each $x \in D$, hence $x^p = x$ for some $n$. It follows that $D$ is commutative by the proposition.

We next generalize the proposition.

Theorem 66.3. (Jacobson) Let $R$ be a rng (i.e., a ring possibly without a one). Suppose for each element $x$ in $R$, there exists an integer $n > 1$ depending on $x$ such that $x^n = x$. Then $R$ is commutative.

Proof. We may assume that $R$ is not the trivial rng. We want to reduce to the proposition. Let $x$ be an element in $R$ and $n > 1$ an integer such that $x^n = x$. Since for every element $z$ in $R$, $z^m = z$ for some integer $m > 1$, $R$ cannot have any nonzero nilpotent elements, i.e., we cannot have $z^m = 0$ unless $z = 0$ – why?

Step 1. The element $x^{n-1}$ is a central idempotent in $R$, i.e., $x^{n-1}$ lies in the center of $R$ and $(x^{n-1})^2 = x^{n-1}$.

Set $e = x^{n-1}$. We have $e^2 = x^{2n-2} = x^n x^{n-2} = xx^{n-2} = x^{n-1} = e$, so $e$ is an idempotent in $R$. If $y$ lies in $R$, then $(ye - ey)^2 = 0 = (ey - ye)^2$.

As $R$ has no nonzero nilpotent elements, we have $ye = ey = ye$, so $e$ is central.

Step 2. Every right ideal $\mathfrak{A}$ in $R$ is a (two-sided) ideal:

Let $r$ be an element in $R$ and $x$ an element in $\mathfrak{A}$. Then, as $\mathfrak{A}$ is a right ideal, with $x^{n-1}$ central and lying in $\mathfrak{A}$, we have $rx = rx^{n-1}x = ((x^{n-1})r)x$ lies in $\mathfrak{A}$.

Step 3. Suppose that $R$ is a ring (i.e., has a one). Define the right radical $\mathfrak{R}$ of $R$ to be the intersection of all maximal right ideals in $R$.

[It can be shown that the left radical of $R$ and the right radical of $R$ are the same, but we do not need this.] Then for all $y$ in $R$, the element $xy - yx$ lies in $\mathfrak{R}$:
Since 1 lies in $R$, the Zorn’s Lemma argument showing that nonzero commutative rings have maximal ideals, shows that maximal right ideals exist in $R$. Let $m$ be any such maximal right ideal. By Step 2, $m$ is an ideal in $R$. By the Correspondence Principle, $R/m$ is a ring with no nontrivial right ideals. It follows that every nonzero element in $R/m$ has a right inverse and therefore an inverse, i.e., $R/m$ is a division ring – why? If $z$ is an element in $R$ with $z^m = z$ and $m > 1$, we have $(z + m)^m = z^m + m = z + m$, so $R/m$ is a division ring satisfying the hypothesis of the proposition, hence commutative. Therefore, $(z + m)(y + m) = (y + m)(z + m)$ for all $y, z$ in $R$, i.e., $zy - yz$ lies in $m$ for any right maximal ideal in $R$ as needed.

**Step 4.** If $R$ is a ring, then $R$ is commutative:

It suffices to show the (right) radical $\mathfrak{R}$ in $R$ is zero. Suppose not and $x$ is a nonzero element in $\mathfrak{R}$. Let $x^n = x$ with $n > 1$ and $e = x^{n-1}$. By Step 1, $e$ is a central idempotent and lies in $\mathfrak{R}$ as $x$ does. Suppose that $1 - e$ is not a unit. Then $1 - e$ lies in a maximal right ideal $m$ in $R$. This implies that 1 lies in $m$ as $e \in \mathfrak{R} \subset m$, which is impossible. Therefore, the right ideal $(1 - e)R = R$, so $(1 - e)r = e$ for some $r$ in $R$. We conclude that $0 = e(1 - e)r = e^2 = e = x^{n-1}$, contradicting $x$ cannot be nilpotent.

**Step 5.** Finish:

By Step 4, we may assume that $R$ does not have a one. Let $e$ be a central idempotent in $R$ as in Step 1 and set $T = eR = Re$. Check that $T$ is a ring with $e = 1$ and for each $y$ in $T$, there exists an integer $n > 1$ such that $y^n = y$. By Step 4, we know that $T$ is a commutative ring. Let $x$ and $y$ be elements of $R$. Then using the commutativity of $T$, we have $xye = (xe)(ye) = (ye)(xe) = yxe$. It follows that $(xy - yx)e = 0$. Since $e$ was a central idempotent in $R$ $(xy - yx)e = 0$ for every central idempotent $e$ in $R$. In particular, if $m > 1$ an integer such that $(xy - yx)^m = xy - yx$ in $R$, then $(xy - yx)^{m-1}$ is a central idempotent by Step 1, so $xy - yx = (xy - yx)(xy - yx)^{m-1} = 0$. This proves that $R$ is commutative.

**Corollary 66.4.** Let $R$ be a ring satisfying for each element $x$ in $R$, there exists an integer $n > 1$ depending on $x$ such that $x^n = x$. Then $R$ is isomorphic to a subring that is a product of fields each algebraic over a finite field of positive characteristic.

**Proof.** We know that $R$ is commutative. Let $R \to \bigtimes_{\text{Min}(R)}^{} R/p$ be the natural map where $\text{Min}(R)$ is the set of minimal prime ideals in $R$. Then each $R/p$ is a commutative domain. Since $x^n = x$ for some $x$, either $x = 0$ or $x^{n-1} = 1$. It follows that $R/p$ is a field and $R$ has
no nonzero nilpotent elements. Therefore, nil(R) = 0 and the map is a ring monomorphism. Finally since each element in R/p, p ∈ Min(R), is a root of a polynomial of the form \(t^m - 1\) over its prime field, we must have R/p must be of characteristic p and be an algebraic extension over its prime field. The result follows \(\square\)

Exercises 66.5.

1. Show that if R is a rng in which every element \(x\) satisfies \(x^n = x\) for some integer \(n\) (depending on \(x\)), that R has no nontrivial nilpotent elements.

2. If R is a nontrivial ring in which every nonzero element has an inverse, show that R is a division ring.

67. Addendum: Hilbert Irreducibility Theorem

Throughout this section, \(x, y, t, t_1, \ldots, t_n\) will be independent variables. Hilbert’s Irreducibility Theorem says if \(f(t_1, \ldots, t_n, t)\) is an irreducible polynomial in \(\mathbb{Q}[t_1, \ldots, t_n, t]\), then \(f(\alpha_1, \ldots, \alpha_n, t)\) is irreducible in \(\mathbb{Q}[t]\) for infinitely many rational numbers \(\alpha_1, \ldots, \alpha_n\). Hilbert used this theorem to prove that there exists a finite Galois extension of \(\mathbb{Q}\) with Galois group \(S_n\) and \(A_n\) for every \(n\). We shall apply the theorem to show \(S_n\) occurs as a Galois group over \(\mathbb{Q}\) for all \(n\). This will show that there exist irreducible polynomials of any degree over \(\mathbb{Q}\) having no root solvable by radicals.

We begin by looking at a polynomial \(f(x, t)\) in \(\mathbb{Q}[x, t]\). Write it as
\[
f(x, t) = a_n(x)t^n + \cdots + a_0(x) \quad \text{with } a_i(x) \in \mathbb{Q}[x] \text{ and } a_n(x) \text{ nonzero,}
\]
so the degree of \(f(x, t)\) in \(t\), \(\deg_t f\), is \(n\). We can view \(f(x, t)\) in \(\mathbb{C}[x, t]\). Since \(\mathbb{Q}\) is infinite, we can and do identify polynomial functions and polynomials over \(\mathbb{C}\). We know by Proposition 41.8 if, in addition, that \(f\) is irreducible in \(\mathbb{C}[x, y]\) that all but finitely many complex numbers \(x_0\) are regular values, i.e., \(a_n(x_0)\) is not zero and \(f(x_0, t)\) has \(n\) distinct roots in \(\mathbb{C}\). We say that a real or complex valued function \(\rho(x)\) is analytic at \(x_0\) if \(\rho\) converges to its Taylor series in a neighborhood of \(x_0\), i.e., has a positive radius of convergence \(R\) centered at \(x_0\). If \(x_0\) is a regular value of \(f(x, t)\), we shall call an analytic function \(\rho = \rho(x)\) a root function of \(f(x, t)\) in a neighborhood of \(x_0\) if \(f(x, \rho(x)) = 0\) in a neighborhood of \(x_0\), i.e., there exists a real number \(R > 0\) such that \(\rho(x)\) are power series in \(x - x_0\) for all \(x\) satisfying \(|x - x_0| < R\). In particular, \(\rho(x_0)\) is a root of \(f(x_0, t)\). We begin by constructing root functions in a neighborhood of regular value \(x_0\) of any given polynomial in \(\mathbb{C}[x, t]\).
Proposition 67.1. Let
\[ f(x, t) = a_n(x)t^n + \cdots + a_0(x) \] with \( a_i(x) \in \mathbb{Q}[x] \)
be a polynomial in \( \mathbb{Q}[x, t] \) of degree \( n \) in \( t \) (so \( a_n(x) \) is nonzero) and \( x_0 \) a regular value of \( f(x, t) \). Then there exist precisely \( n \) distinct root functions \( \rho_1, \ldots, \rho_n \) of \( f(x, t) \) at \( x_0 \) and they satisfy
\[
(*) \quad f(x, t) = a_n(x) \prod_{i=1}^{n} (t - \rho_i(x)) \quad \text{in a neighborhood of } x_0
\]
with the \( \rho_i(x) \) power series expansions in \( x - x_0 \) with coefficients in \( \mathbb{Q} \), the algebraic closure of \( \mathbb{Q} \) in \( \mathbb{C} \), in this neighborhood of \( x_0 \).

**Proof.** Let \( \alpha_1, \ldots, \alpha_n \) be the (distinct) roots of \( f(x_0, t) \) in \( \mathbb{C} \), so lie in a finite extension of \( \mathbb{Q} \), e.g., any field containing \( K = \mathbb{Q}(\alpha_1, \ldots, \alpha_n) \) in \( \mathbb{Q} \). For each \( i = 1, \ldots, n \), we shall construct a unique root function \( \rho_i \) satisfying \( \rho_i(x_0) = \alpha_i \) with coefficients \( K \) in its power series expansion about \( x_0 \). We make the following reductions:

(i) As \( a_n(x_0) \neq 0 \), we may assume that \( a_n(x_0) = 1 \).

(ii) We may assume that \( x_0 = 0 \) by replacing \( f(x, t) \) by the polynomial \( g(x, t) = f(x + x_0, t) \).

(iii) We may assume that \( \alpha_i = 0 \) by replacing \( f(x, t) \) by the polynomial \( g(x, t) = f(x, t - \alpha_i) \).

So we can assume that \( x_0 = 0 \) is a regular value of \( f(x, t) \) having root \( \alpha_i = 0 \) and we must construct a unique root function in a neighborhood of \( x = 0 \) with the given conditions. Therefore, we must find a unique power series
\[
(67.2) \quad \rho(x) = \sum_{i=1}^{\infty} b_i x^i \quad \text{with } b_i \in K \text{ converging for all } |x| < R
\]
for some positive real number \( R \) that satisfies \( f(x, \rho(x)) = 0 \) in a neighborhood of \( x = 0 \). As \( f(0, 0) = 0 \), we can write the polynomial
\[
f(x, t) = a_{10}x + a_{01}t + \sum_{i+j \geq 2} a_{ij}x^i t^j \quad \text{in } \mathbb{Q}[x, t]
\]
with almost all \( a_{ij} \) zero and the \( a_{ij} \in K \). Since \( x_0 = 0 \) is a regular value of \( f(x, t) \), we also have \( \frac{\partial f}{\partial t}(0, 0) = a_{01} \) is nonzero, so we can further assume that \( a_{01}(0, 0) = -1 \), i.e.,
\[
(67.3) \quad f(x, t) = a_{10}x - t + \sum_{i+j \geq 2} a_{ij}x^i t^j \quad \text{in } \mathbb{Q}[x, t].
\]
Assume that $\rho(x)$ has a positive radius of convergence about 0. We show that the $b_i$ in (67.2) are uniquely determined by the equation

\[(67.4)\quad f(x, \sum_{i=1}^{\infty} b_i x^i) = 0\]

in a neighborhood of $x = 0$, i.e., $\rho(x)$ is the unique root function in a neighborhood of 0 and with the $b_i$ lying in $K$.

The coefficient of $x^i$ in the expansion of the left hand side of equation (67.4) must be zero for each $i$. In particular, evaluating the $x$ term shows that $0 = a_{10} - b_1$, so we must have $b_1 = a_{10}$. Inductively, we may suppose that we have constructed unique $b_1, \ldots, b_{k-1}$ in $K$. It is convenient to view $y$ as a new variable to replace $t$. Let

\[h(x, y) = a_{10} x - y + \sum_{i+j \geq 2} a_{ij} x^i y^j \quad \text{and} \quad y_0 = \sum_{i=1}^{k-1} b_i x^i.\]

Then the (formal) Taylor expansion for $h(x, y)$ at $y_0$ exists by Exercise 67.19(2) and is

\[h(x, y) = f(x, y_0) + \frac{\partial f}{\partial y}(x, y_0)(y - y_0) + \text{terms in higher powers of } y - y_0.\]

By (67.3), we must have

\[\frac{\partial f}{\partial y}(x, y) = -1 + g(x, y)\]

with every term in $g(x, y)$ having total degree at least one. Evaluating $y$ at $\rho(x)$ yields $\rho(x) - y_0 = \sum_{i=k}^{\infty} b_i x^i$, so

\[0 = f(x, \sum_{i=1}^{k-1} b_i x^i) + (-1 + g(x, \sum_{i=1}^{k-1} b_i x^i)) (\sum_{i=k}^{\infty} b_i x^i) + \text{terms with power of } x \text{ at least } 2k.\]

As the coefficient of $x^k$ in this equation is zero, the coefficient $-b_k$ of $x^k$ of $f(x, y_0)$ in this equation is uniquely determined by $f$ and $b_1, \ldots, b_{k-1}$ as desired.

We can say more. If $\rho(x)$ is an analytic function, i.e., has positive radius of convergence, then in some neighborhood of $x_0 = 0$, we would have $b_1 = a_{10}$ and $f(x, \sum_{i=1}^{\infty} b_i x^i) = 0$, so

\[\sum_{i=2}^{\infty} b_i x^i = \sum_{i+j \geq 2} a_{ij} x^i (\sum_{i=1}^{\infty} b_i x^i)^j.\]
Expanding this and looking at the coefficient of \(x^k\) for \(k > 1\) leads to an equation

\[ b_k = \sum_{i+j=k} c_{ij}a_{ij}b_l \]

for some nonnegative integers \(c_{ij}\). In particular, \(l < k\). For each of the finite nonzero \(a_{ij}\) occurring in \(f(x, t)\), let \(t_{ij}\) be a new variable. By induction, it follows that there exist nonzero polynomials \(p_k \in \mathbb{Z}[t_{ij}]\) having positive integer coefficients (in the nonzero terms) and satisfying \(b_k = p_k(a_{ij})\).

To finish we must still show that \(\rho(x) = \sum b_kx^k\) has a positive radius of convergence in a neighborhood of \(x_0 = 0\). Choose a positive integer \(A\) satisfying \(A > |a_{ij}|\) for all of the finitely many nonzero \(a_{ij}\) occurring in \(f(x, t)\), and set

\[ f_A(x, t) = Ax - t + A(\sum_{i+j \geq 2} x^it^j). \]

By the argument above, we know if \(f_A\) has a positive radius of convergence, then there exist unique \(\bar{b}_k \in \mathbb{C}\) and \(\bar{p}_k \in \mathbb{Z}[t_{ij}]\) satisfying \(f_A(x, \sum_{i=1}^\infty b_i x^i) = 0\) and \(\bar{b}_k = \bar{p}_k(A, \ldots, A)\) are positive integers for all \(k\). Since \(\bar{b}_k \geq p_k(|a_{ij}|) \geq |b_k|\) for all \(k\), it suffices to show that \(y = \sum_{i=1}^\infty \bar{b}_i x^i\) converges for \(|x| < R\) for some positive real number \(R\).

Summing the geometric series, shows

\[
0 = Ax - y + A \sum_{i+j \geq 2} x^iy^j
= Ax - y + Ax^0 \sum_{j=2}^\infty y^j + Ax^1 \sum_{j=1}^\infty y^j + A \sum_{i=2}^\infty x^i(\sum_{j=0}^\infty y^j)
= Ax - y + \frac{y^2}{1-y} + Ax - \frac{y}{1-y} + A - \frac{x^2}{1-x} - \frac{1}{1-y}.
\]

Multiplying by \(1 - y\) yields

\[
0 = Ax - Ayx - y + y^2 + Ay^2 + Axy + A\frac{x^2}{1-x}
= (A+1)y^2 - y + A\frac{x}{1-x}.
\]

One of the two solutions for \(y(0) = 0\), i.e., when \(x = 0\) is

\[
y = \frac{1 - \sqrt{1 - 4(A+1)Ax/(1-x)}}{2(A+1)}.
\]

Since \(\sqrt{1 - (bx/(1-x))} = (\sqrt{1-cx})/(\sqrt{1-x})\) has a positive radius of convergence about \(x = 0\) for \(b = 4A(A+1)\) and \(c = b+1\), it follows that \(f_A\) is analytic in some neighborhood of \(x_0 = 0\), hence \(\rho(x)\) is a root function in a neighborhood of \(0\).
We therefore conclude that for each \( i = 1, \ldots, n \), there exist a unique root function \( \rho_i(x) \) in a neighborhood of \( x_0 \) for the root \( \alpha_i \) of \( f(x_0, y) \) with coefficients in its power series expansion in \( K \). Since the constant term of \( \rho_i(x) \) is \( \alpha_i \) for \( i = 1, \ldots, n \), the \( \rho_i \) are all distinct in some common neighborhood of \( x_0 \) and satisfy \( f(x, t) = \prod_{i=1}^n (t - \rho_i(x)) \) in this neighborhood. \( \square \)

We shall need the following well-known lemma whose proof we leave as an exercise.

**Lemma 67.5.** (Lagrange Interpolation) Suppose that \( F \) is an infinite field, \( a_0, \ldots, a_n \) distinct elements in \( F \) and \( c_0, \ldots, c_n \) elements in \( F \). Then there exists a unique polynomial \( f \) in \( F[t] \) of degree at most \( n \) satisfying \( f(a_i) = c_i \) for each \( i = 0, \ldots, n \).

Let \( \alpha_0 < \cdots < \alpha_m \) be real numbers. Recall that the \((m+1)\)st Vandermonde determinant is given by

\[
V_m = V_m(\alpha_0, \alpha_1, \ldots, \alpha_m) = \det \begin{pmatrix} 1 & \alpha_0 & \cdots & \alpha_0^m \\ \vdots & \vdots & \ddots & \vdots \\ 1 & \alpha_m & \cdots & \alpha_m^m \end{pmatrix} = \prod_{0 \leq i < j \leq m} (\alpha_j - \alpha_i) \neq 0.
\]

If \( z \) is a real-valued function of \( x \), let

\[
W_m = W_m(\alpha_0, \alpha_1, \ldots, \alpha_m, z) = \det \begin{pmatrix} 1 & \alpha_0 & \cdots & \alpha_0^{m-1} & z(\alpha_0) \\ \vdots & \vdots & \ddots & \vdots & \vdots \\ 1 & \alpha_m & \cdots & \alpha_m^{m-1} & z(\alpha_m) \end{pmatrix}.
\]

With this notation, we have:

**Lemma 67.6.** Let \( \alpha_0 < \cdots < \alpha_m \) be real numbers and \( z(x) \) an \( m \) times differential real-valued function on the closed interval \([\alpha_0, \alpha_m]\). Then there exist a real number \( \alpha \) satisfying

\[
\frac{z^{(m)}(\alpha)}{m!} = \frac{W_m}{V_m}
\]

satisfying \( \alpha_0 < \alpha < \alpha_m \) with \( W_m \) and \( V_m \) as above. In particular, if \( z^{(m)}(x) \) is nonzero on \([\alpha_0, \alpha_m]\), then \( W_m \) is also nonzero.

**Proof.** Let \( y(t) = a_0 + a_1 t + \cdots + a_m t^m \) in \( \mathbb{R}[t] \) be the unique polynomial satisfying \( y(\alpha_i) = z(\alpha_i) \) for \( i = 0, \ldots, m \) given by Lagrange Interpolation. By Cramer’s Rule, the system of equations

\[
a_0 + a_1 \alpha_i + \cdots + a_m \alpha_i^m = z(\alpha_i) \quad \text{for} \quad i = 0, \ldots, m
\]
satisfies \( a_m = W_m/V_m \).

By Rolle’s Theorem there exists \( m \) distinct real numbers \( \alpha_i \) with each \( \alpha_i \) lying in the open interval \((\alpha_i, \alpha_{i+1})\) and satisfying \( y^{(i)}(\alpha_i) = z^{(i)}(\alpha_i) \) for \( i = 0, \ldots, m-1 \). Similarly, we see that there exist \( m - 1 \) distinct points on which \( y^{(2)} \) and \( z^{(2)} \) agree. Inductively, we find that there exists a point \( \alpha \) satisfying
\[
\alpha_0 < \alpha < \alpha_m
\]
on which \( y \) and \( z \) agree. Since \( y^{(m)}(x) = m! a_m \), we have \( z^{(m)}(\alpha) = y^{(m)}(\alpha) = m! a_m \) as needed. \( \square \)

We now prove the Hilbert Irreducibility Theorem in the special case of two variables. We shall use the following observation.

**Observation 67.7.** As \( \text{char}(\mathbb{C}) = 0 \), the domain of polynomials in \( \mathbb{C}[t] \), can be identified with the domain of polynomial functions on an interval \( I \) (of nonzero length) of \( \mathbb{C} \), hence as a subring of the domain \( R \) of convergent functions on \( I \). It follows that we can identify the quotient field of \( \mathbb{C}[t] \) with a subfield of the quotient field of \( R \), the field of finite Laurent series on \( I \). [Note we can also work with polynomials in \( 1/t \) with the appropriate modifications.] We shall utilize this as an identification below.

**Proposition 67.8.** Let \( f(x,t) \) be an irreducible polynomial in \( \mathbb{Q}[x,t] \). Then there exist infinitely many rational numbers \( \alpha \) such that \( f(\alpha, t) \) is irreducible in \( \mathbb{Q}[t] \).

**Proof.** Let 
\[
f(x, t) = a_n(x)t^n + \cdots + a_0(x)
\]
with \( a_i(x) \) in \( \mathbb{Q}[x] \) for \( i = 0, \ldots, n \) and \( a_n \neq 0 \). Let \( d = \max\{\deg_x a_i(x) \mid 0 \leq i \leq n\} \). Almost all elements of \( \mathbb{C} \) are regular values of \( f(x,t) \). Let \( x_0 \) be one such. Replacing \( f(x,t) \) by the irreducible polynomial
\[
h(x,t) = f(x+x_0, t),
\]
we may assume \( x_0 = 0 \). By Proposition 67.1, we can write
\[
(i) \quad f(x, t) = a_n(x) \prod_{i=1}^{n} \left( t - \rho_i(x) \right)
\]
for each \( x \) in some neighborhood of \( x_0 = 0 \) with \( \rho_i \) the \( n \) distinct root functions of \( f \) at \( x = 0 \) (so \( \rho_i(x_0) \) are the roots of \( f(x_0, t) \)) and these root functions have coefficients lying in \( \bar{\mathbb{Q}} \), the algebraic closure of \( \mathbb{Q} \) in \( \mathbb{C} \), in their power series expansion in this neighborhood. It is convenient to work at ‘infinity’ rather than at 0, where a real or complex valued function \( h(x) \) is said to be **analytic at infinity** if it can be represented by a power series \( \sum_{i=1}^{\infty} c_i \frac{1}{x^i} \) converging for all \( |x| > R \) for some positive
real number $R$. Set
\[ g(x, t) = x^d f \left( \frac{1}{x}, t \right), \]
then $g(x, t)$ also lies in $\mathbb{Q}[x, t]$. Moreover, $g(x, t)$ is irreducible. Indeed, suppose that
\[ g(x, t) = r(x, t)s(x, t) \]
with $r(x, t)$ and $s(x, t)$ polynomials in $\mathbb{Q}[x, t]$ of degrees $m_1$ and $m_2$ in $x$, respectively. Then $d = m_1 + m_2$ and
\[ f(x, t) = x^d g \left( \frac{1}{x}, t \right) = x^{m_1} r \left( \frac{1}{x}, t \right) x^{m_2} s \left( \frac{1}{x}, t \right) \]
is a factorization of $f(x, t)$ into polynomials. As $f(x, t)$ is irreducible, this must be a trivial factorization, i.e., either $r$ or $s$ lies in $\mathbb{Q}$ as required.

Now note if $g(\alpha, t)$ is irreducible for $\alpha \in \mathbb{Q}$, so is $f \left( \frac{1}{\alpha}, t \right) = \frac{1}{\alpha^d} g(\alpha, t)$. Therefore, we can replace $f(x, t)$ by $g(x, t)$ and work with power series analytic at infinity, i.e., we may assume all the root functions $\rho_i$ of $f(x, t)$ have power series that converge for all $|x| > R$ and (i) holds for all $|x| > R$ for some real number $R$.

Let $S$ be a nonempty proper subset of $\{1, \ldots, n\}$ and consider the equation
\[ (ii) \quad \prod_{i=1}^{n} (t - \rho_i(x)) = \prod_{i \in S} (t - \rho_i(x)) \prod_{i \notin S} (t - \rho_i(x)). \]
We know that for every value of $x$ with $|x| > R$, we have $f(x, t) = a_n(x) \prod_{i=1}^{n} (t - \rho_i(x))$ is a factorization into linear terms in $\widetilde{\mathbb{Q}}[t]$, since $f(x, t)$ is irreducible in $\mathbb{Q}[x, t]$ and $a_n(x) \in \mathbb{Q}[x]$, neither factor on the right hand side of (ii) can lie in $\mathbb{Q}(x)[t]$ by Lemma 34.7, a consequence of Gauss' Lemma, and Observation 67.7. Consequently,
\[ (iii) \quad \prod_{i \in S} (t - \rho_i(x)) = t^{|S|} + b_1 t^{|S|-1} + \cdots + b_{|S|} \text{ in } \widetilde{\mathbb{Q}}[t], \]
with each $b_j = b_j(x)$, $1 \leq j \leq |S|$, analytic at infinity, but at least one, say $b_i = b_i(x)$ does not lie in $\mathbb{Q}(x)$. Let $N$ be the number of such factorizations (ii), so $N = 2^n - 1$. For each such $S$, let $y_S(x)$ be a coefficient $b_i(x)$ on the right hand side of (iii) analytic at infinity and not in $\mathbb{Q}(x)$.

**Claim.** Let $R' \geq R$ be arbitrary. Then there exists a rational number $x_0$ satisfying $x_0 > R'$ (so a regular value of $f$) satisfying $y_S(x_0)$ is not rational for all nonempty proper subsets $S$ of $\{1, \ldots, n\}$. In particular, $f(x_0, t)$ is irreducible and the proposition is true.
Suppose that \( x_0 \) satisfies the property of the claim, but that \( f(x_0, t) = r(t)s(t) \) in \( \mathbb{Q}[t] \) with \( r \) and \( s \) nonconstant polynomial in \( \mathbb{Q}[t] \). Then there exists a nonempty proper subset \( S \) of \( \{1, \ldots, n\} \) satisfying

\[
    r(t) = a_n(x_0) \prod_{i \in S} (t - \rho_i(x_0)) \quad \text{and} \quad s(t) = \prod_{i \not\in S} (t - \rho_i(x_0))
\]

in \( \mathbb{Q}[t] \). This means that \( y_S(x_0) \) lies in \( \mathbb{Q} \), a contradiction. So we need only construct \( x_0 \) as claimed, since \( R' > R \) is arbitrary.

We construct an \( x_0 \) in \( \mathbb{Q} \) satisfying the claim. In fact, we shall produce an integer \( x_0 \) satisfying the claim. Fix a nonempty proper \( S \) and let \( y = y_S(x) \). Then \( y \) lies in \( \mathbb{Q}[\rho_1(x), \ldots, \rho_n(x)] \subset \mathbb{Q} \), so \( y(x) \) satisfies an equation

\[ (iv) \quad d_l y^l + d_{l-1} y^{l-1} + \cdots + d_0 = 0 \]

with \( d_i \in \mathbb{Q}(x) \) for \( i = 0, \ldots, l \) and \( d_i \) nonzero for some integer \( l \). Clearing denominators, we may assume all the \( d_i \) lie in \( \mathbb{Q}[x] \) and then clearing denominators again that all \( d_i \) lie in \( \mathbb{Z}[x] \).

Multiplying equation \((iv)\) by \( d_l^{-1} \) shows that \( z = z(x) = d_l y(x) \) satisfies an equation

\[ z^l + b_{l-1} z^{l-1} + \cdots + b_0 = 0 \]

with \( b_0, \ldots, b_{l-1} \) all lying in \( \mathbb{Z}[x] \). If \( \alpha > R \) satisfies \( y(\alpha) \) is rational, then \( z(\alpha) \) is an integer as each \( b_i(\alpha) \) is an integer and any rational root of a monic equation with integer coefficients must be an integer. Therefore, we conclude that \( z(x) \) has an integral value at some integer \( \alpha > R \) if and only if \( y(x) \) has a rational value at that \( \alpha \). Consequently, it suffices to work with \( z(x) \). As \( z = d_l y \) with \( d_l \in \mathbb{Z}[x] \) and \( y(x) \) analytic at infinity, we have

\[ (v) \quad z(x) = c_k x^k + \cdots + c_0 + c_{-1} \frac{1}{x} + \text{higher powers of } \frac{1}{x} \]

with \( c_i \in \mathbb{C} \) and \( z \) converges for all \( |x| > R \).

We must show that \( z(\alpha) \) is not an integer for some integer \( \alpha > R' \) for infinitely many integers \( \alpha > R \), hence for any \( R' \geq R \). There are three cases to consider.

**Case 1.** \( z(x) \) is a polynomial in \( \mathbb{C}[x] \):

The polynomial \( z(x) \) can have at most \( k \) integer values by Lagrange Interpolation, lest \( z(x) \) hence \( y(x) \) be a polynomial in \( \mathbb{Q}[x] \), contrary to choice. So \( z(x) \) has no integer values at integer points for sufficiently large \( x \).

**Case 2.** There exists a coefficient \( c_j \) of \( z(x) \) in \((v)\) that is not real:
We may assume that \( j \) has been chosen maximal with \( c_j \) not real. Then we have, where \( \text{Im} \) is the imaginary part,
\[
\lim_{x \to \infty} \text{Im}(\frac{z(x)}{x^j}) = \text{Im}(c_j) \neq 0,
\]
hence \( z(x) \) is not real for all large real values of \( x \).

**Case 3.** All the \( c_j \) in (v) are real and there exist a \( j > 0 \) with \( c_{-j} \) not zero:

Choose \( m = m(S) > 0 \) such that the \( m \)th derivative of \( z(x) \) has no non-negative powers of \( x \), say
\[
z^{(m)}(x) = \frac{d}{x^k} + \text{terms with larger powers of } \frac{1}{x}.
\]
So \( d \) is a nonzero real number and \( k \) a positive integer. In particular, we have
\[
\lim_{x \to \infty} x^k z^{(m)}(x) = d \neq 0.
\]
Choose \( R_1 > R' \) such that if \( \alpha > R_1 \), then
\[
0 < |z^{(m)}(\alpha)| < \frac{2|d|}{\alpha^k}.
\]
Suppose that there exist \( m+1 \) integers \( \alpha_i \) satisfying
\[
R_1 < \alpha_0 < \cdots < \alpha_m \text{ with } z(\alpha_i) \in \mathbb{Z} \text{ for } i = 0, \ldots, m.
\]
Let
\[
V_m = \det \begin{pmatrix} 1 & \alpha_0 & \cdots & \alpha_0^m \\ \vdots & \vdots & \ddots & \vdots \\ 1 & \alpha_m & \cdots & \alpha_m^m \end{pmatrix}
\]
and
\[
W_m = \det \begin{pmatrix} 1 & \alpha_0 & \cdots & \alpha_0^{m-1} & z(\alpha_0) \\ \vdots & \vdots & \ddots & \vdots & \vdots \\ 1 & \alpha_m & \cdots & \alpha_m^{m-1} & z(\alpha_m) \end{pmatrix}.
\]
Applying Lemma 67.6, we see that both \( V_m \) and \( W_m \) are nonzero integers and
\[
\frac{2|d|}{m! \alpha_0^k} \geq \frac{2|d|}{m! \alpha^k} \geq \frac{|z^{(m)}(\alpha)|}{m!} = \frac{|W_m|}{|V_m|} \geq \frac{1}{|V_m|}
\]
for some \( \alpha \) satisfying \( \alpha_0 < \alpha < \alpha_m \). Therefore,
\[
\frac{m!}{2|d|} \alpha_0^k \leq |V_m| = \prod_{i<j} (\alpha_j - \alpha_i) < (\alpha_m - \alpha_0)^{m(m+1)/2}.
\]
Consequently, there exist positive real numbers \( \gamma \) and \( \delta \) depending on \( m \) satisfying
\[
(vi) \quad \alpha_m - \alpha_0 > \gamma \alpha_0^\delta.
\]
Note that the right hand side of \((vi)\) goes to infinity as \(\alpha_0 \to \infty\), i.e., the \(\alpha_m\)’s get further and further from \(\alpha_0\) as it gets larger. Recall that \(N = 2^n - 1\) is the number of the \(y(x)\)’s. (We can replace \(N\) by any larger integer.) Choose \(R_2 > R_1 > R' \ge R\) satisfying
\[
\gamma R_2^\delta > Nm.
\]
In the above, we may also assume that \(\alpha_0 > R_2\). Therefore, we have \(\alpha_0 < \alpha_1 < \cdots < \alpha_m\) with each \(\alpha_i\) and each \(z(\alpha_i)\) integers, \(i = 1, \ldots, m\), and
\[
\alpha_m - \alpha_0 > \gamma \alpha_0^\delta > \gamma R_2^\delta > Nm.
\]
We conclude that among any \(N m + 1\) sufficiently large consecutive integers \(l, l + 1, \ldots, l + Nm\), the function \(z(x)\) takes on integer values on at most \(m\) of these, equivalently \(y(x)\) takes on rational values on at most \(m\) of these.

If \(z(x)\) satisfies Case 1 or 2, it has only finitely many integer values at integer points. Choose \(R_0 > R_2\) such that all such \(z(x)\)’s have no integer values, hence \(y(x)\) no rational values, at integers greater than \(R_0\).

For each of the finite number of \(S\)’s, we have constructed an integer \(m = m(S)\) and a real number \(R_0 = R_0(S) > R'\) such that for any \(x > R_0\) and string of \(N m + 1\) consecutive integers greater than \(R_0\), the analytic function \(y_S(x)\) has at most \(m\) rational values at such a string of \(N m + 1\) consecutive integers. Let
\[
M \ge \max_S \left( m(S) \right) \quad \text{and} \quad R_3 > \max_S \left( R_0(S) \right).
\]
For any string of \(N M + 1\) consecutive integers greater than \(R_3\), each \(y_S(x)\) can have at most \(M\) rational values. In particular, in each string of \(N M + 1\) consecutive integers greater than \(R_3\), there exists an integer \(\alpha\) such that \(y_S(\alpha)\) is not rational for every one of the \(N\) functions \(y_S(x)\). This establishes the claim and completes the proof. \(\square\)

**Corollary 67.9.** Let \(f_1(x, t), \ldots, f_m(x, t)\) be irreducible polynomials in \(Q[x, t]\). Then there exist infinitely many rational numbers \(\alpha\) such that \(f_1(\alpha, t), \ldots, f_m(\alpha, t)\) are all irreducible in \(Q[t]\).

**Proof.** By the proof of the proposition, we can work with the union of all the \(y(x)\)’s that arise from the \(f_i(x, t)\)’s with \(N\) the number of all of these. \(\square\)

To prove the full version of the Hilbert Irreducibility Theorem, we use a trick developed by Kronecker.
Definition 67.10. Let $F$ be a field, $d$ and $n$ positive integers. Set

$$P_d(n, F) := \{f \in F[t_1, \ldots, t_n] \mid \deg_{t_i} f < d \text{ for all } i = 1, \ldots, n\}$$

and

$$S_d : P_d(n, F) \to P_{dn}(1, F)$$

defined by

$$S_d(f(t_1, \ldots, t_n)) = f(t, t^d, t^{d^2}, \ldots, t^{d^{n-1}}).$$

For example, $S_d(t_1^{i_1} \cdots t_n^{i_n}) = t_1^{i_1 + i_2 d + i_3 d^2 + \cdots i_n d^{n-1}}$. Let $m$ be a positive integer, then $m$ has a unique $d$-adic expansion. In particular, applying this to all positive integers $m < d^n$, we deduce that $S_d$ is a bijection. Although $P_d(n, F)$ is not closed under products, we do have:

Remark 67.11. If $f, g, fg \in P_d(n, F)$, then

$$S_d(fg) = S_d(f)S_d(g).$$

In particular, since $S_d$ is a bijection, if $S_d(f)$ is irreducible in $F[t]$, then $f$ is irreducible in $F[t_1, \ldots, t_n]$.

The converse of this remark is false. For example, if $F = \mathbb{Q}$, then $f = t_1^2 + t_2^2$ in $\mathbb{Q}[t_1, t_2]$ is irreducible, but $S_d(f) = t^2 + t^6$ is not. Suppose that we have an irreducible polynomial $f \in P_d(n, F)$ but

$$S_d(f) = G(t)H(t) \text{ in } F[t] \text{ with } 0 < \deg G, \deg H < d^n,$$

i.e., $S_d(f)$ is not irreducible. Then $G$ and $H$ lie in $P_{dn}$, so there exist unique polynomials $g$ and $h$ in $P_d(n, F)$ satisfying $G = S_d(g)$ and $H = S_d(h)$, i.e.,

$$S_d(f) = S_d(g)S_d(h) = S_d(gh).$$

Theorem 67.12. (Kronecker’s Criterion) Let $F$ be a field and $f$ an element in $P_d(n, F)$. Then $f$ is irreducible in $F[t_1, \ldots, t_n]$ if and only if for all non-trivial factorizations (i.e., polynomials of positive degree)

$$S_d(f) = S_d(g)S_d(h)$$

with $g$ and $h$ in $P_d(n, F)$, the polynomial $gh$ contains a nonzero monomial of the form $bt_1^{\ell_1} \cdots t_n^{\ell_n}$ for some $i_j \geq d$.

**Proof.** ($\Leftarrow$): If $f = f_1f_2$ with $f_1, f_2 \in F[t_1, \ldots, t_n]$, then we have $f_1, f_2 \in P_d(n, F)$, $S_d(f) = S_d(f_1f_2) = S_d(f_1)S_d(f_2)$, and $f = f_1f_2$ has no nonzero monomial $bt_1^{\ell_1} \cdots t_n^{\ell_n}$, some $i_j > 0$.

($\Rightarrow$): Suppose that $f$ is irreducible and $S_d(f) = S_d(f_1)S_d(f_2)$ for some $f_1, f_2 \in P_d(n, F)$. If $f_1f_2 \in P_d(n, F)$, then $S_d(f) = S_d(f_1)S_d(f_2) = S_d(f_1f_2)$. As $S_d : P_d(n, F) \to P_{dn}(1, F)$ is one-to-one, $f = f_1f_2$, and one of $f_1, f_2$ lies in $F$. \qed
Example 67.13. If \( f = t_1^2 + t_2^2 \) in \( \mathbb{Q}[t_1, t_2] \), we have \( S_3(f) = t^2(1 + t^4) \), \( S_2(g) = t^2 \) with \( g = t_1^2, \) \( S_3(h) = 1 + t^4 \) with \( h = 1 + t_1t_2, \) and \( gh = t_1^2 + t_2^2. \)

Corollary 67.14. Let \( F \) be a field and \( f \in F[t_0, \ldots, t_n] \) irreducible with \( \deg t_i < d \) for \( i = 1, \ldots, n. \) Let \( S_d : P_d(n, F(t_0)) \rightarrow P_d(1, F(t_0)) \) be defined by \( S_d(f(t_0, t_1, \ldots, t_n)) = f(t_0, x, x^d, \ldots, x^{d^n - 1}). \) Suppose that

\[
(*) \quad S_d(f) = \prod_{i=1}^{m} p_i(t_0, x)
\]

with \( p_i(t_0, x) \in F[t_0][x] \) irreducible for \( i = 1, \ldots, m. \) Then there exists an element \( \varphi \in F(t_0) \) with the following property: If \( \alpha \in F \) satisfies \( \varphi(\alpha) \) is defined and nonzero and \( p_i(\alpha, x) \) is irreducible in \( F[x] \) for \( i = 1, \ldots, n, \) then \( f(\alpha, t_1, \ldots, t_n) \) is irreducible in \( F[t_1, \ldots, t_n]. \)

Proof. As \( F[t_0] \) is a UFD, all factorizations of \( S_d(f) \) into two polynomials in \( F[t_0, x] \) arise from \( (*) \) (up to constants in \( F \)). View \( f \) as an irreducible polynomial in \( F(t_0)[t_1, \ldots, t_n]. \) By Kronecker’s Criterion, any such factorization can also be written \( S_d(f) = S_d(g)S_d(h) \) with \( g, h \in P_d(n, F(t_0)) \) and satisfying some nonzero monomial \( \varphi_{gh}t_1^{i_1} \cdots t_n^{i_n} \) occurs in \( gh \) with \( \varphi_{gh} \in F(t_0) \) (and, in fact, in \( F[t_0] \)), for some \( i_j > d. \)

Let \( \varphi \) be the product of all the \( \varphi_{gh} \) arising from such factorizations. If \( \alpha \in F \) satisfies \( \varphi(\alpha) \) is defined and nonzero and \( p_i(\alpha, x) \) is irreducible in \( F[x] \) for \( 1 \leq i \leq m, \) then \( S_d(f(\alpha, t_1, \ldots, t_n)) = \prod_{i=1}^{m} p_i(\alpha, x) \) is a factorization into irreducibles in \( F[x]. \) It follows that \( f(\alpha, t_1, \ldots, t_n) \) is irreducible in \( F[t_1, \ldots, t_n] \) by Kronecker’s Criterion.

Theorem 67.15. (Hilbert Irreducibility Theorem) Let \( f(x_1, \ldots, x_n, t) \) be an irreducible polynomial in \( \mathbb{Q}[x_1, \ldots, x_n, t]. \) Then there exist infinitely many \( (\alpha_1, \ldots, \alpha_n) \) in \( \mathbb{Q}^n \) such that \( f(\alpha_1, \ldots, \alpha_n, t) \) is irreducible in \( \mathbb{Q}[t]. \)

Proof. We induct on \( n. \) The case of \( n = 1 \) is Proposition 67.8. We shall apply the previous corollary and its notation to complete the proof. If \( f \in P_d(n, \mathbb{Q}(x_1)) \) and \( S_d : P_d(n, \mathbb{Q}(x_1)) \rightarrow P_1(d^n, \mathbb{Q}(x_1)) \), we are done if \( S_d(f) \) is irreducible by induction by choosing \( \alpha \in \mathbb{Q}(x_1) \) such that \( \varphi(\alpha) \) is defined and nonzero. [For almost all \( x_0, \) we have \( \varphi(x_0) \) is defined and nonzero.] So we may assume that we can factor

\[
S_d(f) = \prod_{i=1}^{m} p_i(x_1, t)
\]

as in \( (*) \). By Corollary 67.9, there exist infinitely many \( \alpha \) in \( \mathbb{Q} \) satisfying \( p_i(\alpha, t) \) is irreducible in \( \mathbb{Q}[t] \) for \( i = 1, \ldots, m. \) By the previous corollary and induction, the theorem follows.

It can be shown that the Hilbert Irreducibility Theorem holds with \( F \) replacing \( \mathbb{Q} \) for any number field, i.e., any finite field extension of \( \mathbb{Q} \).
or more generally any non-algebraic finitely generated field extension of any field. However, it does not hold for every field replacing $\mathbb{Q}$, e.g., it does not hold for $\mathbb{C}$ (although it holds for $\mathbb{C}(t)$).

We now apply the Hilbert Irreducibility Theorem to show that $S_n$ is a Galois group over $\mathbb{Q}$.

**Theorem 67.16.** Let $n$ be any positive integer. Then there exists a finite Galois extension $L/\mathbb{Q}$ with $G(L/\mathbb{Q}) \cong S_n$.

**Proof.** Let $K = \mathbb{Q}(t_1, \ldots, t_n)$ and $S_n$ act on $K$ by $\sigma(t_i) = t_{\sigma(i)}$ for $i = 1, \ldots, n$. Set $F = K^{S_n}$. By Example 56.6(3), we know that $F = \mathbb{Q}(s_1, \ldots, s_n)$ with each $s_i$ the $i$th elementary symmetric function of $t_1, \ldots, t_n$, the extension $K/F$ is a finite Galois extension with Galois group $G(K/F) \cong S_n$, and $K$ is a splitting field of

$$f = \prod_{i=1}^{n} (t - t_i) = t^n - s_1 t^{n-1} + \cdots + (-1)^n s_n.$$

By Remark 56.10 to the Primitive Element Theorem, we can find distinct integers $m_1, \ldots, m_n$ satisfying $K = F(m_1 t_1 + \cdots + m_n t_n)$.

Let $c = m_1 t_1 + \cdots + m_n t_n$. Then the $\sigma(c)$, $\sigma \in G(K/F)$ ($\cong S_n$), give $n!$ distinct elements in $K$. Let $g = \prod_{\sigma \in G(K/F)} (t - \sigma(c)) \in K[t]$. Since $\sigma(g) = g$ for all $\sigma \in G(K/F)$, we know that $g \in F[t]$ and $K = F(c)$ is a splitting field of $g$ over $F$. In particular, as $\deg g = [K : F]$ and $m_F(c) | g$ in $F[t]$, we must have $g = m_F(c)$ is irreducible in $F[t]$. Since $s_1, \ldots, s_n$ are algebraically independent, $g(s_1, \ldots, s_n, t)$ is irreducible in $\mathbb{Q}(s_1, \ldots, s_n, t)$, so we can choose $\alpha_1, \ldots, \alpha_n$ in $\mathbb{Q}$, such that $g(\alpha_1, \ldots, \alpha_n, t)$ is irreducible in $\mathbb{Q}[t]$ by the Hilbert Irreducibility Theorem 67.15, i.e., we apply the evaluation map $e_{\alpha_1, \ldots, \alpha_n} : \mathbb{Q}[s_1, \ldots, s_n] \to \mathbb{Q}$ with $s_i \mapsto \alpha_i$ for each $i$. Let $\beta_1, \ldots, \beta_n$ be the roots of the polynomial $f(\alpha_1, \ldots, \alpha_n, t)$ in $\mathbb{Q}[t]$ in its splitting field $L$ in $\mathbb{C}$. Then

$$f(\alpha_1, \ldots, \alpha_n, t) = \prod_{i=1}^{n} (t - \beta_i) = t^n - \alpha_1 t^{n-1} + \cdots + (-1)^n \alpha_n$$

in $\mathbb{Q}[t]$. Let $e = m_1 \beta_1 + \cdots + m_n \beta_n$ in $L = \mathbb{Q}(\beta_1, \ldots, \beta_n)$. As $c = m_1 t_1 + \cdots + m_n t_n$ is a root of $g$ with $t_1, \ldots, t_n$ the roots of $f$, we must have $e$ is a root of the irreducible polynomial $g(\alpha_1, \ldots, \alpha_n, t)$. Therefore,

$$|G(L/\mathbb{Q})| = [L : \mathbb{Q}] \geq \deg m_{\mathbb{Q}}(e) = \deg g(\alpha_1, \ldots, \alpha_n, t) = n!.$$  

It follows that $G(L/\mathbb{Q}) \cong S_n$, since $\deg f = n$ and we can view $G(L/F) \subset S_n$. \qed
Corollary 67.17. There exists an irreducible polynomial of every degree \( n \geq 5 \) not solvable by radicals.

Corollary 67.18. Let \( m \geq 2 \). Then there exists an element algebraic of degree \( 2^m \) over \( \mathbb{Q} \) that is not constructible.

Proof. Let \( f \) be an irreducible polynomial of degree \( n \) in \( \mathbb{Q}[t] \) with Galois group \( S_n \) with \( n = 2^m \) and \( K = F(\alpha) \) a splitting field of \( f \) over \( \mathbb{Q} \). If \( \alpha \) was constructible (from \( z = 0, z = 1 \)), then \( n! = [K : \mathbb{Q}] = 2^e \) for some \( e \) by the refined form of the Constructibility Criterion 59.3, which is impossible.

Exercises 67.19.

1. Let \( F \) be a field of characteristic zero and \( g \) and \( h \) two polynomials in \( F[t] \). Show if the (formal) derivative \( (g - h)' = 0 \), then \( h = g + (h(0) - g(0)) \).

2. Let \( F \) be a field of characteristic zero, \( h \) be a polynomial of degree \( n \) in \( F[t] \) and \( t_0 \) an element of \( F \). Show that the (formal) Taylor expansion of \( h \) exists at \( t_0 \), i.e.,

\[
h(t) = \sum_{i=0}^{n} \frac{h^{(i)}}{i!}(t - t_0)^i
\]

where \( h^{(i)} \) is the \( i \)th derivative of \( h \).

3. Let \( f(x, t) \) be an irreducible polynomial in \( \mathbb{Q}[x, t] \). Show that the set \( \{ \alpha \in \mathbb{Q} \mid f(\alpha, t) \in \mathbb{Q}[t] \text{ is irreducible} \} \) is dense in \( \mathbb{R} \).
Part 6

Additional Topics
CHAPTER XIII

Transcendental Numbers

68. Liouville Numbers

Although we have seen that there are uncountably many real numbers transcendental over \( \mathbb{Q} \), it is usually quite difficult to prove a particular real number that is not easily seen to be algebraic over \( \mathbb{Q} \) is transcendental. In this section, we demonstrate that certain numbers are in fact transcendental over \( \mathbb{Q} \). The first numbers shown to be transcendental over \( \mathbb{Q} \) were constructed by Liouville. His approach depends on the following theorem:

**Theorem 68.1.** (Liouville) Let \( \alpha \) be a complex number algebraic over \( \mathbb{Q} \) of degree \( n > 1 \). Then there exists a positive constant \( c = c(\alpha) \) satisfying

\[
|\alpha - \frac{p}{q}| > \frac{c}{|q|^n}
\]

for all integers \( p \) and \( q \) with \( q \) nonzero.

**Proof.** By definition, \( \deg m_\mathbb{Q}(\alpha) = n \). Clearing denominators, we see that there exists a nonzero integer \( a \) such that \( am_\mathbb{Q}(\alpha) \) lies in \( \mathbb{Z}[t] \). Dividing \( am_\mathbb{Q}(\alpha) \) by its content, we see that there exists an irreducible polynomial \( f \) in \( \mathbb{Z}[t] \) of degree \( n \) with \( f(\alpha) = 0 \). As \( m_\mathbb{Q}(\alpha) \) has no rational roots, neither does \( f \). (Of course, the ideals \( (m_\mathbb{Q}(\alpha)) \) and \( (f) \) are the same in \( \mathbb{Q}[t] \).) As the characteristic of \( \mathbb{Q} \) is zero,

\[
f'(\alpha) = \lim_{z \to \alpha} \frac{f(z) - f(\alpha)}{z - \alpha} = \lim_{z \to \alpha} \frac{f(z)}{z - \alpha}
\]

is nonzero, equivalently

\[
\lim_{z \to \alpha} \frac{|f(z)|}{|z - \alpha|} = 1.
\]

Choose \( \delta > 0 \) such that if \( 0 < |z - \alpha| \leq \delta \), then

\[
|\frac{f(z)}{z - \alpha}| < 2|f'(\alpha)|, \quad \text{so} \quad |f(z)| < 2|f'(\alpha)||z - \alpha|.
\]
Let $p$ and $q$ be integers with $q$ nonzero. Suppose that $0 < \frac{|p|}{|q^n|} \leq \delta$ (as $\alpha$ is not rational), then

$$|f(\frac{p}{q})| < 2|f'(\alpha)||\frac{p}{q} - \alpha|.$$  

Since $\deg f = n$ and $f(\frac{p}{q})$ is nonzero, $f(\frac{p}{q}) = \frac{b}{q^m}$ for some nonzero integer $b$. It follows, if $|\frac{p}{q} - \alpha| \leq \delta$, then

$$\frac{1}{|q^n|} \leq |f(\frac{p}{q})| < 2|f'(\alpha)||\frac{p}{q} - \alpha|,$$

i.e.,

$$\frac{p}{q} - \alpha > \frac{1}{2|f'(\alpha)|} \frac{1}{|q^n|}.$$  

Now suppose that $|\frac{p}{q} - \alpha| > \delta$. Then $\delta \geq \frac{\delta}{|q^n|}$, so $|\frac{p}{q} - \alpha| > \frac{\delta}{|q^n|}$. Set $c := \min\{\delta, \frac{1}{2|f'(\alpha)|}\}$, then $|\frac{p}{q} - \alpha| > \frac{c}{|q^n|}$ for all nonzero integers $q$. \hfill $\Box$

Let $\alpha \in \mathbb{C}$ be a nonrational number. We say that $\alpha$ is a Liouville number if there exist no pair of real numbers $c > 0$ and $n \geq 2$ satisfying

$$0 < |\frac{p}{q} - \alpha| < \frac{c}{|q^n|}$$

for all integers $p$ and $q$ with $q$ nonzero.

Note that if $\alpha$ is a Liouville number, it must be real, since a nonreal complex number has positive distance from any rational number.

**Corollary 68.2.** Liouville numbers are transcendental over $\mathbb{Q}$.

**Remark 68.3.** Let $\alpha$ be a Liouville number, $N$ a positive integer. Then there exist infinitely many rational numbers $\frac{p}{q}$, with $p, q$ with $q \neq 0$ satisfying

$$0 < |\frac{p}{q} - \alpha| < \frac{1}{|q^N|}$$

for all integers $p$ and $q$ with $q$ nonzero.

Indeed suppose this is false, then we could choose a real number $c > 0$ smaller than the distance of $\alpha$ to each of these finitely many rational numbers such that

$$0 < |\frac{p}{q} - \alpha| < \frac{c}{|q^N|}$$

holds for no integers $p$ and $q$ with $q$ nonzero.

We use this to give an alternative characterization of Liouville numbers.

**Proposition 68.4.** Let $\alpha$ be a real number. Then $\alpha$ is a Liouville number if and only if for every positive integer $N$ there exist integers $p$ and $q$, with $q \geq 2$ satisfying

$$0 < |\frac{p}{q} - \alpha| < \frac{1}{|q^N|}.$$  

(*)
Proof. \((\Rightarrow)\) is the remark above.

\((\Leftarrow)\): Let \(c > 0\) and \(n \geq 2\) be given. Choose an integer \(m\) so that \(\frac{1}{2^m} < c\) and set \(N = n + m\). By assumption, there exist integers \(p\) and \(q\), with \(q \geq 2\) satisfying

\[
\left| \frac{p}{q} - \alpha \right| < \frac{1}{q^N} \leq \frac{1}{2^m q^n} < \frac{c}{q^m}.
\]

By the theorem, it follows that \(\alpha\) is a Liouville number if we can show that \(\alpha\) is not a rational number. Suppose to the contrary that \(\alpha = \frac{r}{s}\), with \(r\) and \(s\) integers, \(s\) positive. We may assume that we have chosen \(N\) so that it satisfies \(2^{N-1} > s\). So by assumption, there exist integers \(p\) and \(q\), with \(q \geq 2\) with satisfying (*) However, we also have

\[
\left| \frac{p}{q} - \alpha \right| = \left| \frac{p}{q} - \frac{r}{s} \right| = \left| \frac{ps - qr}{qs} \right| > \frac{1}{sq} > \frac{1}{2^{N-1}q} \geq \frac{1}{q^N},
\]

a contradiction. \(\square\)

Example 68.5. Let \(\alpha = \sum_{k=1}^{\infty} \frac{1}{10^k!}\). Then \(\alpha\) is a Liouville number:

Let \(\sum_{k=1}^{N} \frac{1}{10^k!} = \frac{P_N}{10^N!}\), so \(P_N\) and \(10^N!\) are positive integers satisfying

\[
0 < |\alpha - \frac{P_N}{10^N!}| = \sum_{k=N+1}^{\infty} \frac{1}{10^k} \leq \frac{1}{10^{(N+1)}!} \sum_{k=0}^{\infty} \frac{1}{10^k} \\
\leq \frac{10}{9} \left( \frac{1}{10^{N+1}} \right) < \frac{1}{(10^N)!^N}.
\]

The basic reason that the example is transcendental over \(\mathbb{Q}\) is that the series converges very rapidly. Liouville numbers have this property. Of course, most transcendental numbers do not, e.g., \(\pi\). Mahler showed that \(|\pi - \frac{p}{q}| > q^{-42}\) if \(q \geq 2\).

Much work was done to strengthen Liouville’s theorem. The best theorem was established by Roth who proved the following

Theorem 68.6. (Roth’s Theorem) Let \(\alpha\) be a complex number algebraic over the rational numbers but not rational. Suppose that there exists a non-negative real number \(\mu\) satisfying

\[
|\alpha - \frac{p}{q}| < \frac{1}{q^\mu}
\]

for infinitely many \(p/q\) with \(p, q\) relatively prime integers \(q > 0\). Then \(\mu \leq 2\). Moreover, this is best possible, i.e., if \(\alpha\) is algebraic over \(\mathbb{Q}\) and
\[ \varepsilon > 0, \text{ then} \]
\[ |\alpha - \frac{p}{q}| > \frac{1}{q^{2+\varepsilon}} \]
holds for all but finitely many \( p/q \) with \( p, q \) relatively prime integers \( q > 0 \), and there exists a positive constant \( c = c(\alpha, \varepsilon) \) such that
\[ |\alpha - \frac{p}{q}| > \frac{c}{|q|^{2+\varepsilon}} \]
for all integers \( p \) and \( q \) with \( q \) nonzero,
i.e., is independent of \( n > 0 \).

Roth won the Fields Medal for this work. This subject is part of what is called diophantine approximation. The proof is delicate, but does not use a lot of theory. A proof can be found in Leveque’s book Topics in Number Theory, Volume II, Chapter 4.

Exercises 68.7.
1. Let \( b_0 < b_1 < b_2 < \cdots \) an increasing sequence of positive integers satisfying \( \limsup_{k \to \infty} b_{k+1}/b_k = \infty \) and let \( \{e_k\} \) be a bounded sequence of positive integers. If \( a \geq 2 \) is an integer, show
\[ \sum_{k=0}^{\infty} (-1)^k \frac{a^{b_k}}{b_k} \quad \text{and} \quad \sum_{k=0}^{\infty} \frac{e_k}{a^{b_k}} \]
are Liouville numbers.
2. For any choice of signs, show that
\[ 1 \pm \frac{1}{2^1} \pm \frac{1}{2^2} \pm \frac{1}{2^3} \pm \cdots \]
is a Liouville number.

69. Transcendence of \( e \)

In this section we prove that the real number \( e \) is transcendental over \( \mathbb{Q} \). The ideas give the basis for showing that \( \pi \) is transcendental, which we do in the following section. We give Hermite’s proof.

Theorem 69.1. The real number \( e \) is transcendental over \( \mathbb{Q} \).

Proof. Suppose that there exist integers \( a_0, \ldots, a_m \) not all zero satisfying
\[ a_m e^m + a_{m-1} e^{m-1} + \cdots + a_1 e + a_0 = 0, \]
i.e., \( e \) is a root of the nonzero polynomial \( \sum_{i=0}^{m} a_i t^i \) in \( \mathbb{Z}[t] \), equivalently, \( e \) is algebraic over \( \mathbb{Q} \), by clearing denominators. We may assume that
$a_m$ is nonzero and without loss of generality that $a_0$ is nonzero as well. Define a polynomial function
$$f(x) := \frac{x^{p-1}(x-1)^p(x-2)^p \cdots (x-m)^p}{(p-1)!}$$
and set
$$F(x) := f(x) + f'(x) + f^{(2)}(x) + \cdots + f^{(mp+p-1)}(x)$$
where $p > 2$ is a prime number to be specified later and $f^{(j)}$ is the $j$th derivative of $f$.

Note that $\deg f = mp + p - 1$, so the $(mp + p)$th derivative of $f$, $f^{(mp+p)}(x) = 0$. If $0 < x < m$, then we have
$$|f(x)| < \frac{m^{p-1}m^p \cdots m^p}{(p-1)!} = \frac{m^{mp+p-1}}{(p-1)!}.$$ 

**Check 69.3.** We have the following:

1. $$\frac{d}{dx}(e^{-x}F(x)) = e^{-x}(F'(x) - F(x)) = -e^{-x}f(x)$$

2. $$a_j \int_0^x e^{-x}f(x)dx = a_j [- e^{-x}F(x)]_{x=0}^{x=j} = a_jF(0) - a_je^{-j}F(j).$$

Multiplying equation (2) by $e^j$ and adding yields
$$A := \sum_{j=0}^m a_je^j \int_0^j e^{-x}f(x)dx = \sum_{j=0}^m a_je^jF(0) - \sum_{j=0}^m a_jF(j)$$
$$= -\sum_{j=0}^m a_jF(j) = -\sum_{j=0}^m \sum_{i=0}^{mp+p-1} a_jf^{(i)}(j).$$

**Claim:** $f^{(i)}(j)$ is an integer for all $i$ and $j$ and $p \mid f^{(i)}(j)$ in $\mathbb{Z}$ if $(i,j) \neq (p-1,0)$.

The claim will follow from the following:

**Lemma 69.5.** Let $h(x) = \frac{x^n g(x)}{n!}$ with $g$ a polynomial in $\mathbb{Z}[t]$. Then

1. The number $h^{(j)}(0)$ is an integer for all $j$.
2. The integer $n + 1$ satisfies $n + 1 \mid h^{(j)}(0)$ in $\mathbb{Z}$ if $j \neq n$.
3. If $g(0) = 0$, then $n + 1 \mid h^{(n)}(0)$ in $\mathbb{Z}$.

**Proof.** Let $t^ng(t) = \sum_{j=0}^\infty c_j t^j$ in $\mathbb{Z}[t]$ (almost all $c_j = 0$). So we have $c_j$ are integers and zero if $j < n$. Since $h(x) = \frac{x^n g(x)}{n!}$, we see that $h^{(j)}(0) = c_j j! / n!$. If $j < n$, then $c_j = 0$. Hence $h^{(j)}(0) = 0$ and if
j > n, then \( h^{(j)}(0) = c_j (j! / n!) \) in \( \mathbb{Z} \). If \( j = n \), then \( h^{(n)}(0) = c_n \) in \( \mathbb{Z} \). Finally, if \( g \) has a zero constant term, i.e., \( g(0) = 0 \), then \( c_n = 0 \) and the lemma follows.

We return to the proof of the theorem. We apply the lemma to \( f(x), f(x+1), \ldots, f(x+m) \), respectively to see that \( f^{(i)}(j) \) is an integer for all \( i \) and \( j \) and \( p \mid f^{(i)}(j) \) with the possible exception when \( i = p-1 \) and \( j = 0 \). Therefore, as \( a_i \) are integers, we have

\[
A = - \sum_{j=0}^{m} \sum_{i=0}^{mp+p-1} a_j f^{(i)}(j) \equiv -a_0 f^{(p-1)}(0) \mod p.
\]

We also have \( f^{(p-1)}(0) = (-1)^p (-2)^p \cdots (-m)^p \). (This is just \( g(0) \), if in the lemma, \( g(x) = (x-1)^p \cdots (x-m)^p \).) In particular, if \( p > m \), then \( p \mid f^{(p-1)}(0) \). Consequently, as \( a_0 \neq 0 \), we must have: if \( p \) is a prime satisfying \( p > |a_0| \) and \( p > m \), then

\[
A \equiv -a_0 f^{(p-1)}(0) \not\equiv 0 \mod p.
\]

In particular, \( |A| \geq 1 \). By (69.2) and (69.4), we have

\[
|A| = \left| \sum_{j=0}^{m} a_j e^j \int_0^j e^{-x} f(x) dx \right| \leq \sum_{j=0}^{m} \left| a_j e^j \int_0^j f(x) dx \right|
\]

\[
\leq \sum_{j=0}^{m} |a_j| e^j j^{mp+p-1} \frac{m}{(p-1)!} \leq \sum_{j=0}^{m} |a_j| e^m m^{mp+p-1} \frac{m}{(p-1)!}
\]

\[
= \left( \sum_{j=0}^{m} |a_j| \right) e^m m^{mp+p} \frac{m}{(p-1)!} = \left( \sum_{j=0}^{m} |a_j| \right) e^m (m^{m+1})^p \frac{m}{(p-1)!}.
\]

Choosing \( p \) sufficiently large leads to \( |A| < 1 \), a contradiction. \( \square \)

**Exercise 69.6.** Verify the two statements in Check 69.3.

### 70. Symmetric Polynomials

We discussed symmetric polynomials in the study of Galois Theory. In order to prove that \( \pi \) is transcendental, we shall need a finer discussion. This material is very important in its own right and is fundamental in the study of polynomials as well as its various generalizations.

We first establish some nomenclature for polynomials in finitely many variables over a commutative ring. Let \( R \) be a commutative ring and \( f \) a nonzero polynomial in \( R[t_1, \ldots, t_n] \). Write

\[
f = \sum_{i_1, \ldots, i_n} a_{i_1, \ldots, i_n} t_1^{i_1} \cdots t_n^{i_n}.
\]
Each nonzero term \( a_{i_1,\ldots,i_n} t_1^{i_1} \cdots t_n^{i_n} \) in \( f \) is called a monomial of total degree \( i_1 + \cdots + i_n \) of \( f \) in \( t_1,\ldots,t_n \) and the total degree of \( f \) in \( t_1,\ldots,t_n \) is defined to be the maximum of the total degree of the nonzero monomials of \( f \) in \( t_1,\ldots,t_n \). For each positive integer \( k \), order the subset

\[
\{ a t_1^{i_1} \cdots t_n^{i_n} \mid a \text{ nonzero in } R \text{ with } k = i_1 + \cdots + i_n \}
\]

of \( R[t_1,\ldots,t_n] \) lexicographically, i.e., if \( a, b \in R \) are nonzero and \( k = i_1 + \cdots + i_n = j_1 + \cdots + j_n \), then

\[
 a t_1^{i_1} \cdots t_n^{i_n} \leq b t_1^{j_1} \cdots t_n^{j_n}
\]

if

\[
i_i = j_i \text{ for all } l \text{ or }
\]

there exists an \( l \) satisfying \( i_s = j_s \) for all \( s \leq l \) and \( i_{l+1} < j_{l+1} \).

For example,

\[
t_1 t_2 t_3^2 < t_1^3 t_2 t_3 \text{ and } t_2^2 < t_1 t_2^27
\]

in \( R[t_1,t_2,t_3] \) where \( < \) of course means \( \leq \) but not \( = \). We define the leading term of a nonzero \( f \) in \( R[t_1,\ldots,t_n] \) to be the nonzero monomial of \( f \) that has maximal total degree relative to the lexicographic ordering on \( R[t_1,\ldots,t_n] \). In particular, the total degree of the leading term is the same as the total degree of \( f \), although there can be many nonzero monomials of \( f \) of the same total degree. The coefficient of the leading term is called the leading coefficient.

**Definition 70.1.** Let \( R \) be a commutative ring. A polynomial \( f \) in \( R[t_1,\ldots,t_n] \) is called symmetric in \( t_1,\ldots,t_n \) if for all permutations \( \sigma \in S_n \), we have \( f(t_1,\ldots,t_n) = f(t_{\sigma(1)},\ldots,t_{\sigma(n)}) \).

**Example 70.2.** The following are symmetric in \( R[t_1,\ldots,t_n] \):

1. Any element of \( R \).
2. \( t_1^r + \cdots + t_n^r \) for every positive integer \( r \).
3. \( t_1 t_2 + t_1 t_3 + \cdots + t_1 t_n + t_2 t_3 + \cdots + t_{n-1} t_n = \sum_{i<j} t_i t_j \)
4. \( \sum_{1 \leq i_1 < i_2 < \cdots < i_r \leq n} t_{i_1}^s \cdots t_{i_r}^s \) for each \( 1 \leq r \leq n \) and each positive integer \( s \).

The \( r \)th elementary symmetric polynomials in \( t_1,\ldots,t_n \) is defined to be

\[
s_r(t_1,\ldots,t_n) = \sum_{1 \leq i_1 < i_2 < \cdots < i_r \leq n} t_{i_1} \cdots t_{i_r}.
\]

For example,

\[
s_3(t_1,t_2,t_3,t_4) = t_1 t_2 t_3 + t_1 t_2 t_4 + t_1 t_3 t_4 + t_2 t_3 t_4 \text{ and } s_4(t_1,t_2,t_3,t_4) = t_1 t_2 t_3 t_4.
\]
We shall write \( s_r \) for \( s_r(t_1, \ldots, t_n) \) if \( n \) is clear and we set \( s_0 = 1 \).

**Observation 70.3.** Let \( R \) be a commutative ring, \( S = R[t_1, \ldots, t_n] \), and \( s_i(t_1, \ldots, t_n) \) in \( S \). Then in \( S[t] \), we have

\[
\begin{align*}
 f & := \prod (t + t_i) := t^n + s_1t^{n-1} + \cdots + s_n \\
 g & := \prod (t - t_i) := t^n - s_1t^{n-1} + \cdots + (-1)^n s_n
\end{align*}
\]

with \(-t_1, \ldots, -t_n \) (respectively, \( t_1, \ldots, t_n \)) the roots of \( f \) (respectively, of \( g \)) in \( S \). Set \( S_0 = R[s_1, \ldots, s_n] \subset S = R[t_1, \ldots, t_n] \). Then the subset \( S_0^{S_n} := \{ f \in S \mid f \text{ symmetric in } t_1, \ldots, t_n \} \) satisfies \( S_0 \subset S_0^{S_n} \subset S \). We shall show that \( S_0 = S_0^{S_n} \).

**Theorem 70.4.** (Fundamental Theorem of Symmetric Polynomials)

Let \( R \) be a commutative ring, \( s_i = s_i(t_1, \ldots, t_n) \) in \( R[t_1, \ldots, t_n] \). Then \( R[s_1, \ldots, s_n] = R[t_1, \ldots, t_n]^{S_n} := \{ f \in S \mid F \text{ symmetric in } t_1, \ldots, t_n \} \).

More specifically, let \( f \in R[t_1, \ldots, t_n] \) be symmetric in \( t_1, \ldots, t_n \) of total degree \( k \) in \( t_1, \ldots, t_n \). Then there exists a unique polynomial \( g \) in \( R[s_1, \ldots, s_n] \) satisfying \( f = g \) and the total degree of \( g \) in \( s_1, \ldots, s_n \) is at most \( k \).

**Proof.** **Existence:** Let \( f \in R[t_1, \ldots, t_n] \) be symmetric in \( t_1, \ldots, t_n \) of total degree \( k \) and leading term \( at_1^{i_1} \cdots t_n^{i_n} \) in the lexicographical ordering. In particular, \( k = i_1 + \cdots + i_n \). As \( f \) is symmetric, \( at_\sigma^{i_1} \cdots t_\sigma^{i_n} \) (respectively, \( at_\sigma^{i_1} \cdots t_\sigma^{i_n} \)) is also a nonzero monomial in \( f \) for every \( \sigma \) in \( S_n \). It follows that \( i_1 \geq i_2 \geq \cdots \geq i_n \). Set

\[
j_l := i_l - i_{l-1} \leq i_1 \text{ for } l = 1, \ldots, n
\]

where \( i_{n+1} = 0 \). Therefore, \( j_l \geq 0 \) and \( i_l = j_1 + j_2 + \cdots + j_n \) for every \( l \).

The monomial \( as_1^{j_1} \cdots s_n^{j_n} \) has total degree in \( t_1, \ldots, t_n \) equal to

\[
j_1 + 2j_2 + \cdots + nj_n = i_1 + i_2 + \cdots + i_n = k
\]

and total degree in \( s_1, \ldots, s_n \) equal to

\[
j_1 + j_2 + \cdots + j_n \leq i_1 + i_2 + \cdots + i_n = k.
\]

The leading term of \( as_1^{j_1} \cdots s_n^{j_n} \) as a polynomial in \( t_1, \ldots, t_n \) is

\[
at_1^{j_1}(t_1t_2)^{j_2}(t_1t_2t_3)^{j_3} \cdots (t_1 \cdots t_n)^{j_n} = at_1^{i_1} \cdots t_n^{i_n},
\]

so \( f_1 := f - as_1^{j_1} \cdots s_n^{j_n} \) has leading term in \( t_1, \ldots, t_n \) less than \( as_1^{j_1} \cdots s_n^{j_n} \) in the lexicographic ordering. Iterating the process rids us of all monomials of total degree \( k \) in \( t_1, \ldots, t_n \) except for monomials in \( s_1, \ldots, s_n \) of total degree at most \( k \) in \( s_1, \ldots, s_n \) that when viewed as polynomials
in $t_1, \ldots, t_n$ have total degree $k$. Induction on the total degree of $f$ produces the existence of the required $g$.

**Uniqueness:** Suppose that $f$ has two such expressions. Then we would have an equation

$$\sum_{t_1, \ldots, t_n} a_{i_1, \ldots, i_n} s_1^{i_1} \cdots s_n^{i_n} = 0$$

not all $a_{i_1, \ldots, i_n}$ zero. Let $as_1^{j_1} \cdots s_n^{j_n}$ be the leading term of this expression in $s_1, \ldots, s_n$. Writing this term in $t_1, \ldots, t_n$ would produce a leading term in $t_1, \ldots, t_n$

$$at_1^{i_1} \cdots t_n^{i_n}$$

with $j_l = i_l - i_{l+1}$, i.e., $i_l = j_l + \cdots + j_n$

just as before. As the $t_i$ are indeterminants, we would have $a = 0$, a contradiction. \hfill \square

We deduce the immediate specialization of this result.

**Corollary 70.5.** Let $R$ be a commutative ring and $f \in R[t_1, \ldots, t_n]$ a symmetric polynomial. If $\alpha_1, \ldots, \alpha_n$ are elements in $R$ and $s_i$ the elementary symmetric polynomials in $t_1, \ldots, t_n$, then the evaluation $e_{\alpha_1, \ldots, \alpha_n}$ of $f, f(\alpha_1, \ldots, \alpha_n), \in R[s_1(\alpha_1, \ldots, \alpha_n), \ldots, s_n(\alpha_1, \ldots, \alpha_n)]$.

We shall need the following consequence of this.

**Corollary 70.6.** Let $F$ be a field, $K/F$ a field extension, and $f$ in $F[t]$ a polynomial of degree $n$ with roots $\alpha_1, \ldots, \alpha_n$ in $K$. Suppose that $P \in F[t_1, \ldots, t_n]$ is symmetric in $t_1, \ldots, t_n$. Then $P(\alpha_1, \ldots, \alpha_n)$ lies in $F$.

**Proof.** Using Observation 70.3, we know if $f = a_n t^n + \cdots + a_0$ in $F[t]$, then $\pm s_i(\alpha_1, \ldots, \alpha_n) = a_i/a_n$ lies in $F$. The result follows. \hfill \square

## 71. Transcendence of $\pi$

The proof of the transcendence of $\pi$ over the rationals takes the approach of the analogous statement for $e$, but it is significantly harder. In particular, we shall need the Fundamental Theorem of Symmetric Polynomials.

We begin with the observation that a function $f : \mathbb{Z} \to \mathbb{Z}$ satisfying $f(p) \to 0$ as $p \to \infty$ must satisfy $f(p) = 0$ for large enough primes, since $|f(p)| < 1$ for all primes large enough.

We give Hilbert’s proof of the following famous result.

**Theorem 71.1.** (Lindemann) The real number $\pi$ is transcendental over $\mathbb{Q}$. 

PROOF. Suppose that \( \pi \) is algebraic over \( \mathbb{Q} \). Then \( \alpha_1 := \sqrt{-1} \pi \) is also algebraic over \( \mathbb{Q} \). We have

(i) (Euler) \[ e^{\alpha_1} + 1 = 0. \]

Let \( g_1 = m_{\mathbb{Q}}(\alpha_1) \) in \( \mathbb{Q}[t] \). Suppose that \( \deg g_1 = n \) and \( \alpha_1, \ldots, \alpha_n \) are all the (distinct) roots of \( g_1 \) in \( \mathbb{C} \).

Then by (i), we have

(ii) \[ \prod_{i=1}^{n} (e^{\alpha_i} + 1) = 0. \]

Expanding (ii) yields

(iii) \[ \left( \sum_{1 \leq r \leq n} \sum_{1 \leq i_1 < \cdots < i_r \leq n} e^{\alpha_{i_1} + \cdots + \alpha_{i_r}} \right) + 1 = 0. \]

Let \( s_j(t_1, \ldots, t_n) := \sum_{1 \leq i_1 < \cdots < i_j \leq n} t_{i_1} \cdots t_{i_j} \) be the \( j \)th elementary polynomial in \( \mathbb{C}[t_1, \ldots, t_n] \). Then the monic polynomial

\[ \prod_{i=1}^{n} (t - \alpha_i) = t^n - s_1(\alpha_1, \ldots, \alpha_n)t^{n-1} + \cdots + (-1)^n s_n(\alpha_1, \ldots, \alpha_n) \]

in \( \mathbb{C}[t] \) has precisely \( \alpha_1, \ldots, \alpha_n \) as roots. In particular, this polynomial is just \( g_1 \). It follows that each \( s_j(\alpha_1, \ldots, \alpha_n) \) lies in \( \mathbb{Q} \).

Fix \( r \) with \( 1 \leq r \leq n \) and let \( T_{i_1, \ldots, i_r}, \ 1 \leq i_1 < \cdots < i_r \leq n, \)

be \( k := \binom{n}{r} \) independent variables, ordered lexicographically on the \( i_1, \ldots, i_r \). Let \( s'_j, 1 \leq j \leq k, \) be the elementary symmetric polynomials in the \( T_{i_1, \ldots, i_r} \) evaluated at the corresponding \( \alpha_{i_1} + \cdots + \alpha_{i_r} \). Set \( g_r := t^k - s'_1 t^{k-1} + \cdots + (-1)^k s'_k, \) a polynomial in \( \mathbb{C}[t] \).

Therefore, \( g_r \) has roots

\( \alpha_{i_1} + \cdots + \alpha_{i_r} \) for \( 1 \leq i_1 < \cdots < i_r \leq n. \)

We can view the evaluation of the elementary symmetric polynomials in the \( T_{i_1, \ldots, i_r} \) in two steps:

(a) Evaluate these elementary symmetric polynomials in the \( T_{i_1, \ldots, i_r} \) at \( t_{i_1} + \cdots + t_{i_r} \) for each \( 1 \leq i_1 < \cdots < i_r \leq n. \)

(b) Evaluate the resulting polynomials in (a) at the corresponding \( \alpha_{i_1} + \cdots + \alpha_{i_r} \).

Check 71.2. The polynomials in (a) above are symmetric in \( t_1, \ldots, t_n \).
It follows by Corollary 70.6, the corollary to the Fundamental Theorem of Symmetric Polynomials 70.4, that the polynomials in \((a)\) lie in \(Q[s_1(t_1, \ldots, t_n), \ldots, s_n(t_1, \ldots, t_n)]\) and
\[
s'_j \in Q[s_1(\alpha_1, \ldots, \alpha_n), \ldots, s_n(\alpha_1, \ldots, \alpha_n)] \subset Q.
\]
Consequently, \(g_r \in Q[t]\) for \(r = 1, \ldots, n\).

**Claim:** There exists a polynomial \(g\) in \(Z[t]\) having as its roots in \(C\) precisely all the nonzero \(\alpha_{i_1} + \cdots + \alpha_{i_r}, 1 \leq i_1 < \cdots < i_r \leq n, 1 \leq r \leq n\):
\begin{align*}
M := & \text{precisely the number of } (i_1, \ldots, i_r), 1 \leq i_1 < \cdots < i_r \leq n, \\
& 1 \leq r \leq n, \text{ satisfying } \alpha_{i_1} + \cdots + \alpha_{i_r} = 0.
\end{align*}

Therefore, the non-negative integer \(M\) satisfies \(t^M \mid \tilde{g}\), but \(t^{M+1} \not\mid \tilde{g}\) in \(Q[t]\). Let \(\hat{g} := \tilde{g}/t^M\) in \(Q[t]\). Consequently, \(\hat{g}\) has as roots precisely those \(\alpha_{i_1} + \cdots + \alpha_{i_r}, 1 \leq i_1 < \cdots < i_r \leq n, 1 \leq r \leq n\), that are nonzero. Let \(m\) be the least common multiple of all the denominators of the (reduced) coefficients of \(\hat{g}\) (i.e., if \(\hat{g} = \sum a_i t^i\) with \(a_i, b_i\) relatively prime integers \((b_i \neq 0)\) if \(a_i \neq 0\) (and \(b_i = 1\) otherwise), then \(m\) is the least common multiple of the \(b_i\)'s). Set \(g = m\hat{g}\), a polynomial in \(Z[t]\) that satisfies the claim.

Let \(r = \deg g\) and rename the roots of \(g\) by
\[
\beta_1, \ldots, \beta_r \text{ not necessarily distinct.}
\]
By choice, all the \(\beta_i\) are nonzero, so (iii) becomes
\[(iv) \quad e^{\beta_1} + e^{\beta_2} + \cdots + e^{\beta_r} + k = 0, \text{ for some positive integer } k.
\]
Write
\[
g = ct^r + c_1 t^{r-1} + \cdots + c_r, \text{ a polynomial in } Z[t].
\]
As \(\deg g = r\) and \(g(0)\) is not zero, we have
\[
c \text{ and } c_r \text{ are nonzero integers.}
\]
Now fix a prime \(p\) and set
\[s = rp - 1.
\]
Define
\[
f = \frac{e^{st^p-1}g^p(t)}{(p-1)!}, \quad F = f + f' + f^{(2)} + \cdots + f^{(s+p)}
\]
where \(f^{(j)}\) is the \(j\)th derivative of \(f\).
Note that \( \deg f = p - 1 + pr = p + s \), so \( f^{(p+1)} = 0 \).

It follows that

\[
\frac{d}{dt} \left( e^{-t} F(t) \right) = -e^{-t} f(t) \quad \text{so} \quad e^{-t} F(t) - F(0) = - \int_0^t e^{-x} f(x) dx.
\]

Multiply this equation by \( e^t \) and let \( x = \lambda t \) (\( \lambda \) is a new variable), to get

\[
(v) \quad F(t) - e^t F(0) = t \int_0^1 e^{(1-\lambda)t} f(\lambda t) d\lambda.
\]

Setting \( t = \beta_1, \ldots, \beta_r \) (one at a time) into (v) and summing, we get, using (iv),

\[
(vi) \quad \sum_{i=1}^r F(\beta_i) + kF(0) = \sum_{i=1}^r \beta_i \int_0^1 e^{(1-\lambda)\beta_i} f(\beta_i \lambda) d\lambda.
\]

We shall show that this leads to a contradiction by analyzing

\[
A = \sum_{i=1}^r F(\beta_i) + kF(0), \quad \text{the left hand side of (vi), and}
\]

\[
B = \sum_{i=1}^r \beta_i \int_0^1 e^{(1-\lambda)\beta_i} f(\beta_i \lambda) d\lambda, \quad \text{the right hand side of (vi)}.
\]

**Claim:** If \( p >> 0 \), then \( A \) is a non-zero integer.

To prove this, we have to first evaluate \( \sum_{i=1}^r f^{(j)}(\beta_i) \). We show

(a) \quad \sum_{i=1}^r f^{(j)}(\beta_i) \quad \text{is an integer.}

(b) \quad p \mid \sum_{i=1}^r f^{(j)}(\beta_i) \quad \text{in } \mathbb{Z}.

As \( g(\beta_i) = 0 \) for all \( i \), all terms in

\[
(*) \quad \left. \frac{d^j}{dt^j} \right|_{t=\beta_i} \left( \frac{c^x t^{p-1} g^p(t)}{(p-1)!} \right)
\]

having a \( g^l \) term with \( l > 0 \) are zero. In particular, (a) and (b) hold if \( j < p \).

So we may assume that \( j \geq p \).

**Check 71.3.** \( \frac{1}{pe^t} f^{(j)}(t) \) lies in \( \mathbb{Z}[t] \) for \( j \geq p \). (Cf. Example 2.11.)
Therefore, \( \frac{1}{pc^s} \sum_{i=1}^{r} f^{(j)}(t_i) \in \mathbb{Z}[t_1, \ldots, t_r] \) is symmetric in \( t_1, \ldots, t_r \) of
total degree \( (p-1) + rp - j = s + p - j \leq s \) in \( t_1, \ldots, t_r \) as \( j \geq p \). By
the Fundamental Theorem of Symmetric Polynomials 70.4,

\[ \frac{1}{pc^s} \sum_{i=1}^{r} f^{(j)}(t_i) \text{ lies in } \mathbb{Z}[s_1(t_1, \ldots, t_r), \ldots, s_r(t_1, \ldots, t_r)] \]

and is of total degree in \( s_1(t_1, \ldots, t_r), \ldots, s_r(t_1, \ldots, t_r) \) at most \((p - 1 + pr) - j = s + p - j \leq s\).

We conclude by (†) that the sum

\[ \frac{1}{pc^s} \sum_{i=1}^{r} f^{(j)}(\beta_i) \text{ lies in } \mathbb{Z}[\beta_1, \ldots, \beta_r, t_1, \ldots, t_r] \]

by evaluating the \( t_i \) at the \( \beta_i \). Since

\[ \frac{g}{c} = t^r + \frac{c_1}{c} t^{r-1} + \cdots + \frac{c_r}{c} \]

has roots \( \beta_1, \ldots, \beta_r, \) we see that

\[ s_j(\beta_1, \ldots, \beta_r) = (-1)^j \frac{c_j}{c} \text{ for } j = 1, \ldots, r, \]

so

\[ \frac{1}{pc^s} \sum_{i=1}^{r} f^{(j)}(\beta_i) \text{ lies in } \mathbb{Z}[\frac{c_1}{c}, \ldots, \frac{c_r}{c}] . \]

Since (†) has total degree in \( s_1(t_1, \ldots, t_r), \ldots, s_r(t_1, \ldots, t_r) \) at most \( s \),
it follows that

\[ \frac{1}{p} \sum_{i=1}^{r} f^{(j)}(\beta_i) \text{ lies in } \mathbb{Z}[c_1, \ldots, c_r] \subset \mathbb{Z} , \]

i.e., \( \sum_{i=1}^{r} f^{(j)}(\beta_i) \) is an element of \( p\mathbb{Z} \). This shows (a) and (b).

**Check 71.4.** We have

\[ f^{(j)}(0) = \begin{cases} 0 & \text{if } j \leq p - 2 \\ c^s c_p^p & \text{if } j = p - 1 \\ l_j p & \text{if } j \geq p \text{ some integer } l_j . \end{cases} \]

We conclude that there exists an integer \( N \) satisfying

\[ \left( \sum_{i=1}^{r} F(\beta_i) \right) + kF(0) = N p + k c^s c_p^p \text{ in } \mathbb{Z} . \]

As \( k, c, \) and \( c_r \) are not zero, if \( p \) is chosen such that \( p > \max \{k, |c|, |c_r| \} \),
then \( p \mid \left( \sum_{i=1}^{r} F(\beta_i) \right) + kF(0) \). This shows that \( A \) is a nonzero integer
for all \( p >> 0 \).
We turn to computing $B$.

**Claim:** $B = 0$ for all $p > 0$.

Of course, if we show this, then we have finished the proof. For each $j$ with $1 \leq j \leq r$, let

$$m(j) := \sup_{0 \leq \lambda \leq 1} |g(\beta_j \lambda)|.$$

Hence if $0 \leq \lambda \leq 1$, we have

$$|f(\beta_j \lambda)| \leq \frac{|e^s| \beta_j |^{p-1} m(j)^p}{(p-1)!}.$$

Let

$$N = \max_j \left( \sup_{0 \leq \lambda \leq 1} |e^{(1-\lambda)\beta_j}| \right),$$

then

$$\left| \sum_{j=1}^{r} \beta_j \int_{0}^{1} e^{(1-\lambda)\beta_j} f(\beta_j \lambda) d\lambda \right| \leq \sum_{j=1}^{r} \frac{|e^s| \beta_j |^{p} m(j)^p}{(p-1)!} N,$$

so the integer

$$\sum_{j=1}^{r} \beta_j \int_{0}^{1} e^{(1-\lambda)\beta_j} f(\beta_j \lambda) d\lambda \to 0 \quad \text{as} \quad p \to \infty.$$

It follows (by the remark at the beginning of the section) that this sum must be zero for all $p > 0$. \qed

**Corollary 71.5.** It is impossible to construct a square from a given circle with the same area using only straight-edge and compass.

We state some other results about transcendental numbers over $\mathbb{Q}$ without proof. A complex number $\alpha$ is called an *algebraic integer* if there exists a monic polynomial $f$ in $\mathbb{Z}[t]$ such that $\alpha$ is a root of $f$. Algebraic Number Theory studies the set

$$\mathcal{A} := \{ \alpha \in \mathbb{C} \mid \alpha \text{ is an algebraic integer} \}.$$

This set is, in fact, a domain. If $K/\mathbb{Q}$ is a finite extension of fields, one also studies $\mathcal{A}_K := \mathcal{A} \cap K$, the *ring of algebraic integers* in $K$. The transcendence of $e$ and $\pi$ are special cases of a more general result that Lindemann sketched and Weierstrass rigorously proved, viz.,

**Theorem 71.6.** Suppose that $\alpha_1, \ldots, \alpha_n$ are distinct algebraic integers and $\beta_1, \ldots, \beta_n$ nonzero algebraic integers. Then

$$\beta_1 e^{\alpha_1} + \cdots + \beta_n e^{\alpha_n} \text{ is not zero.}$$
It follows from this theorem that if the $\alpha_1, \ldots, \alpha_n$ are linearly independent over $\mathbb{Q}$, then the $e^{\alpha_1}, \ldots, e^{\alpha_n}$ are algebraically independent over $\mathbb{Q}$. In particular $e$ and $\pi$ are transcendental over $\mathbb{Q}$. In addition, it implies, for example, that $\cos \alpha$, $\sin \alpha$, and $\tan \alpha$ are transcendental over $\mathbb{Q}$ if $\alpha$ is a nonzero algebraic integer and $\log \alpha$ is transcendental over $\mathbb{Q}$ if $\alpha$ is an algebraic integer different from zero or one.

One of the famous Hilbert Problems is:

**Problem 71.7.** (Hilbert Seventh Problem) Let $\alpha$ and $\beta$ be algebraic integers with $\alpha$ not zero or one and $\beta$ not rational, is it true that $\alpha^\beta$ is transcendental over $\mathbb{Q}$, e.g., $\sqrt{2}^{\sqrt{2}}$ is transcendental over $\mathbb{Q}$?

This was solved in the affirmative independently by Gelfand and Schneider. They showed

**Theorem 71.8.** Suppose that $\alpha_1, \alpha_2, \beta_1, \beta_2$ are nonzero algebraic integers. If $\log \alpha_1$ and $\log \alpha_2$ are linearly independent over $\mathbb{Q}$, then $\beta_1 \log \alpha_1 + \beta_2 \log \alpha_2$ is not zero.

This theorem implies the solution of the Hilbert Problem. This was generalized by Baker in the 1966, who showed

**Theorem 71.9.** Suppose that $\alpha_1, \ldots, \alpha_n$ are nonzero algebraic integers such that $\log \alpha_1, \ldots, \log \alpha_n$ are linearly independent over $\mathbb{Q}$. Then $1, \log \alpha_1, \ldots, \log \alpha_n$ are linearly independent over the quotient field of $\mathcal{A}$.

**Remark 71.10.** The theorem implies the following: Let $\alpha_1, \ldots, \alpha_n$ be nonzero algebraic integers and $\beta_0, \ldots, \beta_n$ be algebraic integers with $\beta_0 \neq 0$. Then

1. $\beta_0 + \beta_1 \log \alpha_1 + \cdots + \beta_n \log \alpha_n \neq 0$.
2. If, in addition, all of the $\beta_i$ are nonzero, then $e^{\beta_0} \alpha_1^{\beta_1} \cdots \alpha_n^{\beta_n}$ is transcendental over $\mathbb{Q}$.
3. If, in addition, none of the $\alpha_i$ are one and $1, \beta_1, \ldots, \beta_n$ are linearly independent over $\mathbb{Q}$, then $\alpha_1^{\beta_1} \cdots \alpha_n^{\beta_n}$ is transcendental over $\mathbb{Q}$.

Baker won the Fields Medal for this work.

**Exercises 71.11.**
1. Prove Check 71.2.
2. Prove Check 71.4.
3. Prove Check 71.3.
CHAPTER XIV

The Theory of Formally Real Fields

In this chapter, we develop the theory of formally real fields including the standard results from Artin-Schreier Theory. In particular, we determine the general form of the Fundamental Theorem of Algebra and solve Hilbert’s seventeenth problem.

72. Orderings

Definition 72.1. A field $F$ is called formally real if $-1$ is not a sum of squares in $F$, i.e., the polynomial $t_1^2 + \cdots + t_n^2$ has no nontrivial zero over $F$ for any (positive) integer $n$. A formally real field is called real closed if it has no proper algebraic extension that is formally real.

If $F$ is formally real, then the characteristic of $F$ must be zero for if char $F = p > 0$, then $-1$ is a sum of $p - 1$ squares.

Notation 72.2. For a commutative ring $R$, let

$$\sum R^2 := \{x \in R \mid x \text{ is a sum of squares in } R\}.$$ 

A field $F$ of characteristic different from two is not formally real if and only if $F = \sum F^2$. In general, $\sum F^2$ is closed under addition and multiplication while $\sum (F^2)^\times := \sum (F^2) \setminus \{0\}$ is a multiplicative group.

Definition 72.3. Let $R$ be a commutative ring, $P \subset R$ a subset. We say that $P$ is a preordering of $R$ if $P$ satisfies all of the following:

1. $P + P \subset P$.
2. $P \cdot P \subset P$.
3. $-1 \notin P$.
4. $\sum R^2 \subset P$.

Let

$$\mathcal{Y}(R) := \{P \mid P \subset R \text{ is a preordering }\}.$$ 

[Of course, $\mathcal{Y}(R)$ may very well be empty.] A preordering $P \in \mathcal{Y}(R)$ is called an ordering if, in addition,

$$R = P \cup -P \text{ and } P \cap -P = \{0\}$$

where $-P := \{x \mid -x \in P\}$.
Let $P$ be a preordering on $R$. We shall show below that if $P$ is a maximal preordering (rel $\subset$), then $R = P \cup -P$ and $P \cap -P$ is a prime ideal in $R$. In particular, if $R$ is a field, this means that $P$ is a maximal preordering if and only if $P$ is an ordering. Note: We have not excluded 0 from lying in $P$.

Let 

$$\mathcal{X}(R) := \{P \in \mathcal{Y}(R) \mid P \text{ is an ordering of } R\}.$$ 

By definition, a field $F$ is formally real if and only if $\sum F^2$ is a preordering. In general, the preordering $\sum F^2$ in a formally real field $F$ is not an ordering. We shall see that $\sum F^2$ is an ordering if and only if $|\mathcal{X}(F)| = 1$. For example, this is the case when a formally real field $F$ is euclidean, i.e., if every element in $F$ is a square or the negative of a square. For such a field, we even have $\sum F^2 = F^2$. For example, the real numbers or the real constructible numbers are euclidean fields. Of course, in general $F^2 \neq \sum F^2$. When a field of characteristic different from two satisfies $F^2 = \sum F^2$, it is called pythagorean. If $F$ is not formally real (and char $F \neq 2$), then $F$ is pythagorean if and only if it is quadratically closed, i.e., $F = F^2$.

**Lemma 72.4.** Let $R$ be a commutative ring, $P \in \mathcal{Y}(R)$.

1. If $a, b \in R$ satisfy $ab \in P$ then either $P + aP \in \mathcal{Y}(R)$ or $P - bP \in \mathcal{Y}(R)$.
2. If $P \in \mathcal{Y}(R)$ is maximal (relative to $\subset$), then $R = P \cup -P$ and $P \cap -P$ is a prime ideal in $R$. In particular, if $R$ is a field then $P \in \mathcal{X}(R)$.
3. If $R$ is a field, then $P$ is a maximal preordering if and only if $P \in \mathcal{X}(R)$.
4. There exists a $\tilde{P} \in \mathcal{Y}(R)$ containing $P$ and satisfying $R = P \cup -P$ and $P \cap -P$ is a prime ideal in $R$. In particular, if $R$ is a field, there exists an ordering $\tilde{P} \in \mathcal{X}(R)$ containing $P$.

**Proof.** (1): We need only show that $-1$ cannot lie in both $P + aP$ and $P - bP$. If it does then there exist $x, y, z, w \in P$ satisfying $-1 = x + ay = z - bw$, so

$$(ay)(-bw) = (-1 - x)(-1 - z) = 1 + x + z + xz,$$

hence $-1 = x + z + xz + abyw \in P$, a contradiction.

(2): If $a \in R$ then $a^2 \in P$, so by (1), either $P + aP \in \mathcal{Y}(R)$ or $P - aP \in \mathcal{Y}(R)$. By maximality, either $a \in P$ or $-a \in P$, so $R = P \cup -P$. Certainly, $P \cap -P$ is closed under addition. If $x \in P \cap -P$ and $y \in R$ then $xy = (-x)(-y) \in P \cap -P$, so $P \cap -P$ is an ideal. Suppose that $ab = (-a)(-b) \in P \cap -P$. If $a \notin P \cap -P$, we may...
assume that \( a \in -P \setminus P \). Maximality and (1) imply that \( -b \in P \). As \( a(-b) \in P \cap -P \) also, the same argument shows that \( b \in P \).

(3): If \( P \in \mathcal{X}(R) \), then \([R^\times : P \setminus \{0\}] = 2\), so \( P \) must be a maximal preordering.

(4) follows from (2) and Zorn’s Lemma. \( \square \)

**Proposition 72.5.** Let \( F \) be a formally real field and \( P \in \mathcal{Y}(F) \). Then

\[
P = \bigcap_{P \subset \mathcal{P} \in \mathcal{X}(F)} \mathcal{P}.
\]

**Proof.** The inclusion \( P \subset \bigcap_{P \subset \mathcal{P} \in \mathcal{X}(F)} \mathcal{P} \) is clear. Conversely, suppose that \( x \notin P \). By definition, \(-1 \notin P \). We show that \( P - xP \notin \mathcal{Y}(F) \).

To do this, it suffices to show that \(-1 \notin P - xP \). If this is false, write \(-1 = y - xz \) with \( y, z \in P \). Then \( z \neq 0 \) so

\[
x = 1 \cdot z^{-1} + yz^{-1} = \frac{z}{z^2} + \frac{yz}{z^2} \in P,
\]

a contradiction. By the lemma, there exists \( \mathcal{P} \in \mathcal{X}(F) \) such that \( P - xP \subset \mathcal{P} \). As \( x \notin P - xP \), lest \(-x^2 \in P - xP \), it follows that \(-x \in \mathcal{P} \), so \( x \notin \bigcap_{P \subset \mathcal{P} \in \mathcal{X}(F)} \mathcal{P} \). \( \square \)

**Corollary 72.6.** (Artin-Schreier) Suppose that \( F \) is formally real. Then \( \sum F^2 = \bigcap_{\mathcal{X}(F)} P \). In particular, \( \mathcal{X}(F) \neq \emptyset \).

**Proof.** As \( F \) is formally real, \( \sum F^2 \in \mathcal{Y}(F) \). \( \square \)

**Remark 72.7.** Let \( F \) be a formally real field and \( P \in \mathcal{X}(F) \). An element \( 0 \neq x \in F \) is called positive (respectively, negative) rel \( P \) and also written \( x >_P 0 \) (respectively, \( x <_P 0 \)) if \( x \in P \) (respectively, \(-x \in P \)). If \( x \in \bigcap_{\mathcal{X}(F)} Q \), i.e., \( x \) is positive at every ordering of \( F \) then \( x \) is called totally positive (and \(-x \) is called totally negative). In particular, we have

\[
\bigcap_{\mathcal{X}(F)} Q \setminus \{0\} = (\sum F^2)^\times,
\]

i.e., \( 0 \neq x \) is totally positive if and only if it is a sum of squares.

If \( P \in \mathcal{X}(F) \), let \( \leq_P \) on \( F \) denote the total ordering given by \( x \leq_P y \) if \( y - x \in P \). Of course, if \( x \leq_P y \) and \( z \in F \) then \( x + z \leq_P y + z \) and if, in addition, \( 0 \leq_P z \) then \( xz \leq_P yz \). We let \( >_P \) and \( \geq_P \) have the obvious meaning.

**Exercises 72.8.**
1. Prove that the field of real constructible numbers $C$ is euclidean and $C(\sqrt{-1})$ is quadratically closed.

2. Show that a formally real field is euclidean if and only if it is pythagorean and has a unique ordering.

3. Show the intersection of pythagorean fields is pythagorean.

4. Let $F$ be a formally real field and $\tilde{F}$ an algebraic closure of $F$. A pythagorean field $L$ is called a pythagorean closure of $F$ if either $L = F$ or whenever $F \subset K < \tilde{F}$, then $K$ is not pythagorean. Show that if $K$ is a field in $\tilde{F}$ containing $F$ and is the intersection of all pythagorean fields containing $F$ in $\tilde{F}$, then $K$ is pythagorean. In particular, a pythagorean closure of $F$ in $\tilde{F}$ exists and is unique. It is denoted by $F_{pyth}$.

5. Let $F$ be a formally real field and $\tilde{F}$ an algebraic closure. Call a square root tower $F \subset F_0 \subset F_1 \subset \cdots \subset F_n$ admissible if for each $i, 1 \leq i \leq n - 1$, there exist $x_i, y_i$ in $F_{i-1}$ satisfying $F_i = F_{i-1}(\sqrt{x_i^2 + y_i^2})$. Show the union of all admissible square towers over $F$ is the pythagorean closure $F_{pyth}$ of $F$. In particular, if $F$ is a formally real field, then the pythagorean closure of $F$ in $\tilde{F}$ exists.

73. Extensions of Ordered Fields

**Definition 73.1.** Let $K/F$ be a field extension with $K$ formally real. Let $P \in \mathcal{X}(K)$. The pair $(K, P)$ is called an ordered field. If $P_0 \in \mathcal{X}(F)$ satisfies $P_0 = P \cap F$, then $(K, P)/(F, P_0)$ is called an extension of ordered fields and $P$ is called an extension of $P_0$. If there exists no proper algebraic extension $L$ of $F$ with $Q \in \mathcal{X}(L)$ satisfying $Q \cap F = P_0$, we say that $(F, P_0)$ is real closed (rel $P_0$). Let $(K, P)/(F, P_0)$ be an extension of ordered fields. If $(K, P)$ real closed (rel $P$) and $K/F$ is algebraic, $(K, P)$ is called a real closure of $(F, P_0)$ (rel $P_0$).

We shall show that given $P \in \mathcal{X}(F)$, there exists a real closure of $(F, P)$ and it is unique up to a (unique) $F$-isomorphism. We begin with the following useful result characterizing extensions of orderings under certain conditions.

**Theorem 73.2.** Let $(F, P)$ be an ordered field.

1. Let $d \in F$ and $K = F(\sqrt{d})$. Then there exists an extension of $P$ to $K$ if and only if $d \in P$.

2. If $K/F$ is finite of odd degree then there exists an extension of $P$ to $K$. 

Proof. (1): Suppose $d \in P$. We may assume that $d \in P \setminus F^2$. Let

$$S = \{ \sum x_i y_i^2 \mid x_i \in P, \ y_i \in K \text{ for all } i \}.$$ 

We show that $S \in \mathcal{Y}(K)$. Certainly, $S$ is closed under addition and multiplication and contains $\sum K^2$. Thus it suffices to show that $-1 \not\in S$. If $-1 \in S$, we can write $-1 = \sum x_i(a_i + b_i\sqrt{d})^2$ for some $a_i, b_i \in F$ and $x_i \in P$. As $\{1, \sqrt{d}\}$ is an $F$-basis for $K$, we arrive at the contradiction that $-1 = \sum x_i(a_i^2 + b_i^2 d) \in P$. This proves that $S \in \mathcal{Y}(K)$. By Lemma 72.4, there exists $\tilde{P} \in \mathcal{X}(K)$ satisfying $P \subset S \subset \tilde{P}$. As $P$ is an ordering, $P = \tilde{P} \cap F$. Suppose that $d \not\in P$. If $\tilde{P} \in \mathcal{X}(K)$ contains $P$ then $d = (\sqrt{d})^2 \in \tilde{P} \cap F = P$, a contradiction.

(2): As the char $F$ is zero, by the Primitive Element Theorem, $K = F(x)$, for some $x \in K$. Let $f$ be the minimal polynomial of $x$ and suppose that

$$S = \{ \sum a_i y_i^2 \mid a_i \in P, \ y_i \in K \text{ for all } i \} \not\in \mathcal{Y}(K).$$

As $K \cong F[t]/(f)$, we have an equation,

$$-1 \equiv \sum a_i g_i(t)^2 \mod (f)$$

for some $a_i \in P$ and $0 \neq g_i \in F[t]$ with $\deg g_i < \deg f$ for all $i$. Lift this to an equation

$$-1 = \sum a_i g_i(t)^2 + qf$$

in $F[t]$. Since $-1 \not\in P$, evaluating at $x$ shows that there exists an $i$ with $\deg g_i > 0$. Let $m = \max_i(\deg g_i) > 0$. Then the coefficient of $t^{2m}$ in $\sum a_i g_i(t)^2$ cannot be zero as it is of the form $\sum a_m b_i^{2m} \geq 0$. Thus $\deg qf = 2m$. As $\deg f = n$ is odd, so is $\deg q$. By construction, $\deg q < \deg f$. Let $h \mid q$ in $F[t]$ with $h \in F[t]$ irreducible of odd degree. Then $-1 \equiv \sum a_i g_i^2 \mod (h)$, so we may repeat the argument. By induction, this reduces to the case when $f$ is linear where the result is false. $\square$

The first part of the theorem has the immediate corollary.

**Corollary 73.3.** Let $(F, P_0)$ be an ordered field. Suppose that there exists an element $d$ in $P_0 \setminus F^2$. Then there exists an proper extension of ordered fields $(F(\sqrt{d}), P)/(F, P_0)$ with $P_0 + dP_0 \subset P$. In particular $F(\sqrt{d})$ is formally real.

We also conclude that

**Corollary 73.4.** $(F, P)$ is real closed (rel $P$) if and only if $F$ is real closed. Moreover, if this is the case, then $P = F^2$. 
Proof. Suppose that \((F, P)\) is real closed. Let \(d \in P\). By the theorem, \(P\) extends to \(F(\sqrt{d})\) so \(d \in F^2\) by hypothesis. As \((F, P)\) is real closed, we must have \(P = F^2 \in \mathcal{X}(F)\). Since \(F^2 \subset Q\) for all \(Q \in \mathcal{X}(F)\), we have \(\mathcal{X}(F) = \{F^2\}\). Suppose that \(L/F\) is algebraic and \(L\) formally real. Let \(Q \in \mathcal{X}(L)\). Then \((L, Q)/(F, F^2)\) is an extension of ordered fields as \(F^2 \subset Q\), hence \(L = F\) and \(F\) must be real closed.

If \(F\) is real closed, then \(F(\sqrt{d})/F\) cannot be a nontrivial extension of formally real fields. It follows by the Corollary 73.3 and Lemma 72.4 that \(F^2 \in \mathcal{Y}(F)\) is an ordering hence \(\mathcal{X}(F) = \{F^2\}\). If \(L/F\) is a proper algebraic extension, then \(-1 \in L^2\) so \(F \cap L^2\) cannot be a preordering on \(F\). It follows that \((F, F^2)\) is real closed. \(\Box\)

There exist fields such that \(\mathcal{X}(F) = \{F^2\}\) but \(F\) not real closed, i.e., euclidean fields that are not real closed. The real constructible numbers is such an example. The theorem only implies that a euclidean field has no formally real quadratic extensions. It also means that \(F(\sqrt{-1})\) is quadratically closed. Indeed if \(\alpha = a + b\sqrt{-1}\) in \(F(\sqrt{-1})\) with \(a, b \in F\), then there exist \(c, d \in F\) satisfying

\[
c^2 = \frac{a + \Delta}{2} \quad \text{and} \quad d^2 = \frac{-a + \Delta}{2}
\]

where \(\Delta = \sqrt{a^2 + b^2}\) is the positive square root in \(F^2\) in euclidean \(F\).

Proposition 73.5. Every ordered field \((F, P)\) has a real closure \((\widetilde{F}, \widetilde{F}^2)\).

Proof. Let \(\widehat{F}\) be an algebraic closure of \(F\). Let \(P \in \mathcal{X}(F)\). The statement follows from a Zorn’s Lemma argument on

\[\{\widehat{F}/K/F \mid (K, Q)/(F, P)\text{ is an extension}\}\]. \(\Box\)

Because of the last results, if \((K, P)/(F, P_0)\) is a real closure, then \(P = K^2\) and \(K\) is real closed so we just say \(K\) is a real closure of \(F\) rel \(P\).

Corollary 73.6. Let \(F\) be a formally real field and \(P \in \mathcal{X}(F)\). Then there exists a real closure of \((F, P)\).

Theorem 73.7. (Fundamental Theorem of Algebra) The following are equivalent:

1. \(F\) is real closed.
2. \(F\) is euclidean and every polynomial \(f \in F[t]\) of odd degree has a root in \(F\).
3. \(F\) is not algebraically closed but \(F(\sqrt{-1})\) is.
Proof. (1) \(
Rightarrow\) (2): We have seen the hypothesis implies that \(F\) is euclidean. Let \(p \in F[t]\) be irreducible of odd degree. As \(K = F[t]/(p)\) is an extension of odd degree any ordering of \(F\) would extend, hence \(p\) is linear. (2) easily follows.

(2) \(
Rightarrow\) (3): This is the essentially the same proof given for the usual Fundamental Theorem of Algebra 56.12. Let \(K = F(\sqrt{-1})\). Then \(K\) is quadratically closed. Let \(E/K\) be an algebraic extension and \(L/F\) the normal closure of \(K/F\). If \(H\) is a 2-Sylow subgroup of the Galois group \(G\) of \(L/F\), then \(L^H = F(\theta)\) for some \(\theta\) by the Primitive Element Theorem, so we must have \(L^H = F\) by hypothesis. Consequently, \(G = H\) and hence must be trivial as \(K\) is quadratically closed.

(3) \(
Rightarrow\) (1): It suffices to show that \(F\) is formally real, i.e., \(-1 \notin \sum F^2\). As \(\sqrt{-1} \notin F\), it suffices to show \(F\) is pythagorean. Let \(a, b \in F^\times\). We need only show that \(a^2 + b^2 \in F^2\). Let \(f = t^4 - 2at^2 + (a^2 + b^2) \in F[t]\). Then \(f = (t^2 - (a + b\sqrt{-1}))(t^2 - (a - b\sqrt{-1}))\) in \(F(\sqrt{-1})\). As \(F(\sqrt{-1})\) is quadratically closed, there exist \(\alpha, \beta \in F(\sqrt{-1})\) satisfying

\[
\alpha^2 = a + b\sqrt{-1} \quad \text{and} \quad \beta^2 = a - b\sqrt{-1},
\]

hence

\[
f = (t - \alpha)(t + \alpha)(t - \beta)(t + \beta) \quad \text{in} \quad F(\sqrt{-1})[t].
\]

As \(f \in F[t]\) cannot be irreducible and \(ab \neq 0\), either \((t - \alpha)(t - \beta)\) or \((t - \alpha)(t + \beta)\) is an irreducible factor of \(f\) over \(F\). Thus \(\alpha \beta \in F\), hence

\[
a^2 + b^2 = (a + b\sqrt{-1})(a - b\sqrt{-1}) = \alpha^2 \beta^2
\]

lies in \(F^2\). \(\square\)

It follows that over a real closed field, the only irreducible polynomials are linear polynomials and irreducible quadratic polynomials of the form \(at^2 + bt + c\) with negative discriminant (in the unique ordering).

Artin and Schreier proved that the conditions of the Fundamental Theorem of Algebra were also equivalent to: \(F\) is not algebraically closed and its algebraic closure is a finite extension of it. We shall come back to this below. Before showing this, we turn to the question of uniqueness.

We first need some facts and notation about symmetric bilinear forms over a field of characteristic different from two. Details are given Appendix D.

Remark 73.8. Let \(V\) be a finite dimensional vector space over a field \(F\) of characteristic different from two and \(B : V \times V \to F\) a symmetric bilinear form, i.e., linear in each variable and \(B(x, y) = B(y, x)\) for all \(x, y \in V\). We call \((V, B)\) a bilinear space. We say \(x, y\) in \(V\) are
orthogonal (rel $B$) if $B(x, y) = 0$. We call subspaces $W_1, W_2$ of $V$ orthogonal if $B(x, y) = 0$ for all $x \in W_1$, $y \in W_2$ and then denote $W_2 + W_2$ by $W_1 \perp W_2$. Every $(V, B)$ has an orthogonal basis $\mathcal{C}$, i.e., $V = \bigoplus_{v \in \mathcal{C}} Fv$. If $B = \{v_1, \ldots, v_n\}$ is an ordered basis for $V$, then we write $[B]_\mathcal{B}$ for the matrix representation $((B(v_i, v_j))$ of $B$. So $\mathcal{B}$ an orthogonal basis if and only if $[B]_\mathcal{B}$ is diagonal. If $\mathcal{C}$ is another basis, then $[1_V]_{\mathcal{B}, \mathcal{C}}[B]_{\mathcal{B}}[1_V]_{\mathcal{B}, \mathcal{C}} = [B]_\mathcal{C}$. We also have $B(x + y, x + y) = B(x, x) + 2B(x, y) + B(y, y)$, so $B$ is determined by $B(x, x), x \in V$ (or $v$ in a basis for $V$).

Let $A$ be a finite dimensional commutative $F$-algebra. If $a \in A$ let $\lambda_a : A \to A$ be the $F$-linear map given by $x \mapsto ax$. The trace $\text{trace} : A \to F$ is the trace of a matrix representation of $\lambda_a$ relative to a fixed basis. It is independent of the choice of basis. If $A$ is a finite separable field extension of $F$ of characteristic different from two, then this is just $\text{Tr}_{A/F}$ (cf. Exercise 59.26(6).) Define the trace form of $A$

$$\varphi = \varphi_A : A \times A \to F$$

to be the symmetric bilinear form given by $\varphi(x, y) = \text{trace}(\lambda_{xy})$.

Let $B$ be a symmetric bilinear form over a formally real field $F$ and $P \in \mathcal{X}(F)$. As an ordering on $F$ is induced by the real closure relative to this ordering, the results of Appendix D hold. In particular, we have the following: [Cf. Proposition D.12 and Corollary D.16 in Appendix D.]

If $V$ is the underlying vector space for $B$, then $V$ has an orthogonal decomposition $V = V_+ \oplus V_- \oplus V_0$ where if $0 \neq v \in V$ then

$$B(v, v) = \begin{cases} > P 0 & \text{if } v \in V_+ \\ < P 0 & \text{if } v \in V_- \\ = 0 & \text{if } v \in V_0 \end{cases}$$

with the dimensions of $V_0$, $V_+$, and $V_-$ independent of such an orthogonal decomposition. The integer

$$\text{sgn}_P B := \dim V_+ - \dim V_-$$

is an invariant of $B$, called the signature of $B$ at $P$, so $\text{sgn}_P B = \text{sgn}_P [B]_\mathcal{B}$ for any basis $\mathcal{B}$ of $V$.

If $K$ is an extension field of $F$ and $\mathcal{B}$ a basis for $V$, then we let $V_K$ be the vector space over $K$ on basis $\mathcal{B}$. By restricting scalars to $F$, we see it is also a vector space over $F$ and there exists a (natural) linear embedding of $V$ into $V_K$ by sending $\mathcal{B}$ to $\mathcal{B}$. We call this extension of scalars. [Cf. making a real vector space into a complex one by extending scalars.] If $B$ is the bilinear form above, define
$B_K: V_K \times V_K \to K$ to be the map additive in each variable satisfying $B_K(\alpha v, \beta w) = \alpha \beta B_K(v, w)$ for all $\alpha$ and $\beta$ in $K$ and $v, w \in \mathcal{B}$. Then, if $(K, Q)/(F, P)$ is an extension, $\text{sgn}_P B = \text{sgn}_Q(B_K)$.

**Computation 73.9.** Let $(F, P)$ be an ordered field, $f \in F[t] \setminus F$ and $A = F[t]/(f)$. Let $- : F[t] \to A$ be the canonical epimorphism and $\varphi : A \times A \to F$ the trace form.

**Case 1.** $f = (t - a)^n$ in $F[t]$:
Let $\mathcal{B} = \{1, \bar{t} - a, \ldots, (\bar{t} - a)^{n-1}\}$, an $F$-basis for $A$. As $(\bar{t} - a)$ is nilpotent, $\text{trace}(\lambda(\bar{t} - a)^i) = 0$ for all $i \geq 1$. It follows that the matrix representation of $\varphi$ in the $\mathcal{B}$ basis is given by the diagonal matrix

$$[\varphi]_\mathcal{B} = \text{trace}(\lambda(\bar{t} - a)^i(\bar{t} - a)^j)) = \begin{pmatrix} n & 0 & \cdots & 0 \\ 0 & 0 & \ddots & \vdots \\ \vdots & \ddots & \ddots & \vdots \\ 0 & \cdots & 0 & 0 \end{pmatrix}.$$ 

Consequently, $\text{sgn}_P \varphi = \text{sgn}[\varphi]_\mathcal{B} = 1$.

**Case 2.** $f = (t^2 + at + b)^n$ in $F[t]$ with $a^2 - 4b < P$ 0:
Let $\mathcal{B} = \{1, \bar{t}, \bar{t}^2 + 1, \bar{t}(\bar{t}^2 + 1), (\bar{t}^2 + 1)^2, \bar{t}(\bar{t}^2 + 1)^2, \ldots, (\bar{t}^2 + 1)^n, \bar{t}(\bar{t}^2 + 1)^n\}$, an $F$-basis for $A$. Computation shows the matrix representation of $\varphi$ in the $\mathcal{B}$ basis is given by the diagonal matrix

$$[\varphi]_\mathcal{B} = \begin{pmatrix} 2n & 0 & \cdots & 0 \\ 0 & -2n \\ \vdots & 0 & \ddots & \vdots \\ 0 & \cdots & 0 & 0 \end{pmatrix}.$$ 

Consequently, $\text{sgn}_P \varphi = 0$.

The next result, due to Sylvester, is the main ingredient in proving the uniqueness of real closures relative to a fixed ordering. It replaces the more common use of Sturm’s Theorem.

**Lemma 73.10.** (Sylvester) Let $(F, P)$ be an ordered field and $K$ a real closure of $(F, P)$. Suppose that $f$ is a non-constant polynomial in $F[t]$ and $A = F[t]/(f)$. Let $\varphi : A \times A \to F$ be the trace form. Then

$$\text{sgn}_P \varphi = \text{the number of roots of } f \text{ in } K.$$ 

In particular, the number of roots of $f$ in $K$ depends only on $P$ and is independent of the real closure $K$. 

Theorem 73.12. Let \( \varphi \) be the trace form on \( A = F[t]/(f) \) and \( f = \prod \varphi_i^{e_i} \) a factorization of \( f \) into monic irreducibles in \( K[t] \) with \( u \in K^\times \). By the Chinese Remainder Theorem, we have a natural ring (and \( K \)-algebra) isomorphism.

\[
A_K := K[t]/(f) \cong \prod K[t]/(p_i^{e_i}) = \prod A_i,
\]

where \( A_i = K[t]/(p_i^{e_i}) \). We view this as an identification of \( K \)-algebras. It is easy to check that \( \varphi_{A_K}(x_i, x_j) = 0 \) if \( x_i \in A_i \) and \( x_j \in A_j \) with \( i \neq j \) (by writing down an appropriate \( K \)-basis), hence we can restrict \( \varphi_A \) to each \( A_i \) to obtain \( \text{sgn}_P \varphi = \text{sgn}_K^2(\varphi_{A_K}) = \sum_{i=1}^r \text{sgn}_K^2 \varphi_{A_i} \). Since \( K \) is real closed, it follows from the Fundamental Theorem of Algebra that \( p_i \) is either linear or quadratic. Therefore, by the computation, we see that

\[
\text{sgn}_P \varphi = \sum_{i=1}^r \text{sgn}_K^2 \varphi_{A_i} = \sum_{p_i \text{ linear}} e_{ij}.
\]

The result follows. \( \square \)

Lemma 73.11. Let \( (F, P) \) be an ordered field and \( K \) a real closure of \( (F, P) \). Suppose that \( E = F(\alpha) \) is an algebraic extension of \( F \) with \( \varphi = \varphi_E : E \times E \to F \) the trace form. Suppose that \( r = \text{sgn}_P \varphi > 0 \). Then the minimal polynomial \( f \) of \( \alpha \) has \( r \) roots in \( K \). Let \( \alpha_1, \alpha_2, \ldots, \alpha_r \) be these roots. Then there exist at most \( r \) distinct and at least one extension \( \tilde{P}_i \) of \( P \) to \( E \). These are induced by the \( r \) distinct \( F \)-homomorphisms \( \sigma_i : E \to K \) given by \( \alpha \mapsto \alpha_i \) with \( \tilde{P}_i = \sigma_i^{-1}(K^2) \).

Proof. Every \( F \)-embedding of \( E \) into an algebraic closure of \( K \) must take \( \alpha \) to a root of \( f \). By Lemma 73.10, the polynomial \( f \) has \( r \) distinct roots in \( K \). Clearly, \( \tilde{P}_i = \sigma_i^{-1}(K^2) \supset P \). Suppose that \( Q \in \mathcal{X}(E) \) extends \( P \) with \( Q \neq \tilde{P}_i \) for \( 1 \leq i \leq r \). Then there exists an \( a_i \in E \) satisfying \( a_i \in Q \) but \( \sigma_i(a_i) \notin K^2 \) for all \( i \). By repeated application of Theorem 73.2, there exists \( \tilde{Q}_i \in \mathcal{X}(E(\sqrt{a_1}, \ldots, \sqrt{a_r})) \) extending \( Q \). By the Primitive Element Theorem, there exists \( \beta \) such that \( F(\beta) = E(\sqrt{a_1}, \ldots, \sqrt{a_r}) \). Let \( L \) be a real closure of \( (F(\beta), \tilde{Q}) \) hence of \( (F, P) \). The minimal polynomial \( g \in F[t] \) of \( \beta \) has root \( \beta \in L \), hence a root in \( K \) by Lemma 73.10. Therefore, there exists an \( F \)-homomorphism \( \tau : F(\beta) \to K \). As the \( F \)-homomorphism \( \tau|_E : E \to K \) must satisfy \( \tau|_E = \sigma_i \) for some \( i \), we have \( \sigma_i(a_i) = \tau(a_i) = (\tau(\sqrt{a_i}))^2 \) lies in \( K^2 \), a contradiction. \( \square \)

Theorem 73.12. Let \( (F, P) \) be an ordered field.
(1) Let \((E, Q)/(F, P)\) be an algebraic extension of ordered fields and \(K\) a real closure of \((F, P)\). Then there exists an order preserving \(F\)-homomorphism \(E \rightarrow K\).

(2) Let \(K_1\) and \(K_2\) be two real closures of \((F, P)\). Then there exists a unique \(F\)-homomorphism \(K_1 \rightarrow K_2\) and it is an order preserving isomorphism.

(3) In Lemma 73.11 and its notation, all the \(\tilde{P}_i\) are distinct. In particular, there exist precisely \(r\) extensions of \(P\) to \(E\).

**Proof.** (1): Apply Zorn’s Lemma to 
\[\{M \mid E/M/F \text{ and there exists an order preserving } F\text{-homomorphism } \psi_M : (M, Q \cap M) \rightarrow (K, K^2)\}\]
to obtain a maximal such \(M_0\). Suppose \(M_0 \neq E\) and \(x \in E \setminus M_0\). Let \(L\) be a real closure of \((E, Q)\). Then \(f\) has a root in \(L\) so \(\text{sgn}_{Q \cap M_0}(x) > 0\). By Lemma 73.11, we can extend \(\psi_{M_0} : (M_0, Q \cap M_0) \rightarrow (K, K^2)\) to \(\psi_{M_0}(x) : (M_0(x), Q \cap M_0(x)) \rightarrow (K, K^2)\), a contradiction. So \(M_0 = E\).

(2): By (1) there exists an \(F\)-homomorphism \(\sigma : K_1 \rightarrow K_2\). Let \(\tau : K_1 \rightarrow K_2\) be another \(F\)-homomorphism. Let \(\alpha \in K_1\) and \(\alpha_1, \ldots, \alpha_r\) be the (distinct) roots of the minimal polynomial \(f\) of \(\alpha\) in \(K_1\). We may assume that
\[\alpha_1 <_{K_1^2} \cdots <_{K_1^2} \alpha_r\]
in \(K_1\). As \(\sigma(K_1^2) \subset K_2^2\) and \(\tau(K_1^2) \subset K_2^2\) and \(\mathcal{X}(K_i) = \{K_i^2\}\) for \(i = 1, 2\) both \(\sigma\) and \(\tau\) are order preserving. Thus
\[\sigma(\alpha_1) <_{K_2^2} \cdots <_{K_2^2} \sigma(\alpha_r)\text{ and } \tau(\alpha_1) <_{K_2^2} \cdots <_{K_2^2} \tau(\alpha_r)\text{ in } K_2.\]

As these are the roots of \(f\) in \(K_2\), we have \(\sigma(\alpha_i) = \tau(\alpha_i)\) for all \(i\). In particular, \(\sigma(\alpha) = \tau(\alpha)\). Hence \(\sigma = \tau\).

(3): Suppose in Lemma 73.11 that we have \(\tilde{P}_i = \tilde{P}_j\) with \(i \neq j\). Let \(\tilde{P} = \tilde{P}_i\) and \(M\) a real closure of \((E, \tilde{P})\). By (1), the embeddings \(\sigma_i \neq \sigma_j\) can be extended to \(F\)-homomorphisms \(M \rightarrow K\). This contradicts (2). \(\square\)

**Exercise 73.13.** Show that the field of real constructible numbers \(C\) has algebraic extensions that are not euclidean. In particular, \(C(\sqrt{-1})\) is not algebraically closed.

**74. Characterization of Real Closed Fields**

We turn to the proof of the generalization of the Fundamental Theorem of Algebra for fields. [We shall prove a version of the usual statement of the Fundamental Theorem of Algebra for division rings in
Theorem 105.12.] The proof is rather intricate in order to take care of the positive characteristic case. The proof revolves around the study of primitive $p^t$th roots of unity for a prime $p$. This is necessary in order to reduce to $\sqrt{-1}$, a primitive fourth root of unity in characteristic zero.

To prove the last Artin-Schreier characterization of real closed fields, we need a few facts from field theory. These are

**Facts 74.1.** Let $F$ be a field.

1. Let $p$ be a prime. Then $\mathbb{Z}/p^r\mathbb{Z}$ is cyclic unless $p = 2$ in which case $\mathbb{Z}/2^r\mathbb{Z}$ is not cyclic if $r \geq 3$. (Cf. Section E.)
2. If $\text{char } F = p > 0$ and $a \in F \setminus F^p$, then $t^p - a \in F[t]$ is irreducible for all $e \geq 1$. (Cf. Example 49.15(8).)
3. If $\text{char } F = p > 0$ and $a \in F$ satisfies $a \neq u^p - u$ for any $u \in F$, then $t^p - t - a \in F[t]$ is irreducible. (Cf. Exercise 55.22(9).)
4. If $\text{char } F = p > 0$ and $a \in F$ satisfies $a \neq u^p - u$ for any $u \in F$ and $E$ is a splitting field of $t^p - t - a$ over $F$, then there exists an extension $K/E$ of degree $p$.
5. Let $p$ be a prime different from $\text{char } F$. If $t^p - 1$ splits in $F[t]$ and $K/F$ is a cyclic extension of degree $p$, then $K = F(r)$ for some $r$ satisfying $r^p \in F$. (Cf. Theorem 59.20.)
6. Suppose $p$ is a prime different from $\text{char } F$ and $F$ contains a primitive $p$th root of unity. If $E/F$ is cyclic of degree $p$, then there exists an $x \in E$ such that $E = F(x)$ and the minimal polynomial of $x$ over $F$ is $t^p - a \in F[t]$. (Same as the previous.)
7. If $\text{char } F = p > 0$ and $E/F$ is cyclic of degree $p$, then there exists an $x \in E$ such that $E = F(x)$ and $x^p - x \in F$. (Same as previous.)

We prove Fact 74.1(4), leaving the others as exercises.

**Proof.** (of Fact 74.1(4)) Let $E = F(u)$ with $m_F(u) = t^p - t - a$, a cyclic extension of degree $p$ with $F$-basis $B = \{1, u, \ldots, u^{p-1}\}$. We show that the element $au^{p-1}$ is not of the form $v^p - v$, $v \in E$. If this is not so, say $au^{p-1} = v^p - v$. Write $v = v_0 + v_1 u + \cdots + v_{p-1} u^{p-1}$ for some $v_0, \ldots, v_{p-1}$ in $E$. Since $u^p = u + a$, we have

$$au^{p-1} = v^p - v = v_0^p + v_1^p (u + a) + v_2^p (u + a)^2 + \cdots + v_{p-1}^p (u + a)^{p-1}$$

$$- v_0 - v_1 u - \cdots - v_{p-1} u^{p-1}.$$

As $B$ is a basis, we see that $v_{p-1}^p - v_{p-1} = a$, a contradiction. By Fact (3), we know that $t^p - t - au^{p-1} \in E[t]$ is irreducible, hence its splitting field is of degree $p$ over $E$. \qed
Theorem 74.2. (Artin-Schreier) Let $C$ be algebraically closed and $R$ a proper subfield with $C/R$ finite. Then $R$ is real closed and $C = R(\sqrt{-1})$.

**Proof.** Let $C' = R(\sqrt{-1})$. As $C$ is the algebraic closure of $C'$, it suffices to show $C = C'$. We have:

$$\text{(74.3)} \quad \text{If } C/E/C', \text{ then } [E : C'] \leq [C : C'] < \infty.$$

**Claim 74.4.** $C'$ is perfect:

Suppose not, then $\text{char } F = p > 0$ and there exists an element $a \in C' \setminus (C')^p$. By Fact 74.1(2), for each $n \geq 1$, there exists $E_n/C'$ of degree $p^n$ contradicting (74.3). Therefore $C'$ is perfect.

By Claim 74.4, we conclude that $C/C'$ is finite Galois, since $C$ is algebraically closed. Let $G = G(C/C')$. We may assume that $C' < C$. In particular, there exists a prime number $p$ dividing $|G|$ and a cyclic subgroup $H \subset G$ of order $p$. Let $E$ be the fixed field $C^H$, so $C/E$ is cyclic of degree $p$.

**Claim 74.5.** $E$ contains a primitive $p$th root of unity and $p \neq \text{char } E$:

Suppose to the contrary that $p = \text{char } E$. By Fact 74.1(7), there exists an $x \in C$ such that $C = E(x)$ is the splitting field of the irreducible polynomial $t^p - t - a \in E[t]$ over $E$. Hence $a = x^p - x$. By Fact 74.1(4), there exists a field extension $K/C$ of degree $p$ which is impossible.

Let $\text{char } E = l \geq 0$. We have shown that $p \neq l$. As $t^p - 1 = (t - 1)(t^{p-1} + \cdots + 1)$ splits in $C[t]$ and each irreducible factor has degree at most $p - 1$, the polynomial $t^p - 1$ must split in $E[t]$, hence $E$ contains a primitive $p$th root of unity. As $\text{char } E \neq p$, we have $C = E(r)$ where the irreducible polynomial of $r$ is $t^p - a \in E[t]$ by Fact 74.1(5). Let $\zeta \in C$ be a primitive $p^2$-root of unity. As $t^{p^2} - a$ splits in $C$, we can write

$$t^{p^2} - a = \prod_{i=0}^{p^2-1} (t - \zeta^i s) \quad \text{with } a = s^{p^2}$$

in $C[t]$.

**Claim 74.6.** There exists a primitive $p^2$-root of unity $\varepsilon$ in $C$ such that $C = E(\varepsilon)$:

If $\zeta^i s \in E$ for some $i$, then $(\zeta^i s)^p \in E$. This would mean that $(\zeta^i s)^p = s^{p^2} = a$ and $t^p - a$ has a root in $E$ contradicting the irreducibility of $t^p - a \in E[t]$. Consequently, $\zeta^i s \notin E$ for all $i$. This means that every irreducible factor of $t^{p^2} - a \in E[t]$ must have degree
p. Let \( f \) be such an irreducible factor, say the minimal polynomial of \( \zeta \)'s. The constant term of \( f \) must be \( b = \zeta^n s^p \in E \) for some \( n \). As \( s^p \not\in E \), we have

\[
C = E(s^p) = E(b\zeta^{-n}) = E(\zeta^n).
\]

Since \( E \) contains a primitive \( p \)th root of unity, \( \varepsilon = \zeta^n \) must be a primitive \( p^2 \)-root of unity as needed.

Let \( \Delta \) be the prime subfield in \( E \).

**Claim 74.7.** There exists an \( r \) such that \( \Delta(\varepsilon) \) contains a primitive \( p^r \)th root of unity but not a \( p^{r+1} \)th primitive root of unity:

Let \( \zeta_{p^r} \) be a primitive \( p^r \)th root of unity in \( C \). Suppose that \( \Delta \cong \mathbb{Q} \). Then \( [\mathbb{Q}(\zeta_{p^r}) : \mathbb{Q}] = p^{r-1}(p - 1) \to \infty \) as \( r \to \infty \). As \( [\mathbb{Q}(\varepsilon) : \mathbb{Q}] < \infty \), Claim 74.7 holds in this case. Suppose that \( \Delta \cong \mathbb{Z}/IZ \). By Claim 74.5, we know that \( l \neq p \). As \( t^{p^r} - 1 \) has no repeated roots in \( C \), we have \( |\Delta(\zeta_{p^r})| \geq p^r \to \infty \). As \( |\Delta(\varepsilon)| < \infty \), Claim 74.7 also holds if \( \Delta \) has positive characteristic.

Let \( r \) be as in Claim 74.7. As \( \varepsilon \in F(\varepsilon) \) and \( \varepsilon \) is a primitive \( p^2 \)-root of unity, we have \( r \geq 2 \). Let \( \omega \) be a primitive \( p^{r+1} \)th root of unity in the algebraically closed field \( C \). Let \( N \) be the Galois group of \( \Delta(\omega)/\Delta \). If \( \Delta \) is a finite field, then \( N \) is cyclic. If \( \mathrm{char} \Delta = 0 \), then \( N \cong (\mathbb{Z}/p^{r+1}\mathbb{Z})^\times \), which is cyclic unless \( p = 2 \) and \( r \geq 3 \) by Fact 1. Moreover, if \( N \) is cyclic, it contains a unique subgroup of order \( p \), hence \( \Delta(\omega) \) contains a unique subfield over which it has degree \( p \).

**Claim 74.8.** \( \Delta(\omega) \) contains two subfields over which it has degree \( p \).

In particular, \( p = 2 \) and \( \mathrm{char} E = 0 \):

Let \( f \in E[t] \) be the minimal polynomial of \( \omega \). Since \( \varepsilon \not\in E \) also \( \omega \not\in E \). Hence \( C = E(\omega) \) and \( \deg f = p \). We also have \( f \mid t^{p^{r+1}} - 1 \in E[t] \) and \( t^{p^{r+1}} - 1 = \prod_{i=0}^{p^r-1}(t - \omega^i) \) in \( C[t] \). Let \( K = \Delta(\omega) \cap E \). Then by the above, \( f \in K[t] \) so \( f \) is the minimal polynomial of \( \omega \) over \( K \). Consequently, \( [\Delta(\omega) : K] = p \) and \( K \) is one of the desired subfields.

Let \( z = \omega^p \), a primitive \( p^r \)th root of unity and \( K' = \Delta(z) \). As \( \Delta(\omega) \) contains a primitive \( p^r \)th root of unity and hence all such, we have \( K' = \Delta(z) \subset \Delta(\omega) \). We show that \( [\Delta(\omega) : \Delta(z)] = [\Delta(\omega) : K'] = p \). As \( \omega \) is a root of \( t^p - z \in K'[t] \), it suffices to show \( t^p - z \in K'[t] \) is irreducible. By Fact 74.1(5), it suffices to show that \( t^p - z \) has no roots in \( K' \). Since \( z \in K' \), certainly \( K' \) contains all the primitive \( p \)th roots of unity. In particular, if \( t^p - z \) has a root in \( K' \), it splits in \( K'[t] \) and \( \omega \in K' \). But this means that \( \Delta(\omega) = K' \subset \Delta(\varepsilon) \) contradicting Claim 74.7. Thus \( [\Delta(\omega) : K'] = p \).
To prove Claim 74.8, we need only show that $K \neq K'$. Suppose that $K = K'$. Then $z \in K' = K$ means that $K$ contains a primitive $p^r$th root of unity. As $K = \Delta(\omega) \cap E$, we have $z \in E$, i.e., $E$ contains a primitive $p^r$th root of unity for some $r \geq 2$, hence it contains a $p^2$-root of unity. Thus $\varepsilon \in E$ which is a contradiction. We conclude that $K \neq K'$ and Claim 74.8 is established.

Thus we have shown that $\text{char } \Delta = 0$ and $C = E(\varepsilon)$ with $\varepsilon$ a primitive $2^2$ root of unity, i.e, $\varepsilon = \pm \sqrt{-1}$. But $\sqrt{-1} \in E$ by Claim 74.5, a contradiction and the theorem is proved. 

Let $F$ be a field and $C$ an algebraic closure. Then

$$F_{\text{sep}} := \{x \mid x \in C \text{ separable over } x\},$$

is called the separable closure of $F$. It is the maximal separable field extension of $F$ in $C$. The Galois group $G(F_{\text{sep}}/F)$ is called the absolute Galois group of $F$.

**Corollary 74.9.** Let $F$ be a field. If the absolute Galois group of $F$ contains a nontrivial element of finite order, then $F$ is formally real and the element is an involution.

**Proof.** Let $C$ be an algebraic closure of $F$ and $F_{\text{sep}}$ the separable closure of $F$ in $C$. Let $\sigma$ lie in $G(F_{\text{sep}}/F)$ be of order $r > 1$. Then $G(F_{\text{sep}}/F_{<\sigma>})$ is finite with $[F_{\text{sep}} : F_{<\sigma>}] = r$. By Theorem 74.2 the result follows if $F$ is perfect, so we may assume that $\text{char } F = p > 0$. Suppose that $C/K/F_{\text{sep}}$ is an intermediate field and $\sigma : K \to C$ an embedding. Then $\sigma$ induces a unique isomorphism $K[t]/(t^{p^n} - a) \to \sigma(K)[t]/(\sigma(t^{p^n} - \sigma(a)))$ taking the unique root of $t^{p^n} - a$ in $K[t]/(t^{p^n} - a)$ in $C$ to the corresponding root of $t^{p^n} - \sigma(a)$ in $\sigma(K)[t]/(t^{p^n} - \sigma(a))$ for all $n$ and all $a \in K$. It follows that if $\sigma : F_{\text{sep}} \to F_{\text{sep}}$ lifts to a unique $F$-automorphism $\tilde{\sigma}$ of $C$ of order $r$. Then $C/C_{<\tilde{\sigma}>}$ is a finite extension of order $r > 1$. The result follows by Theorem 74.2.

**Corollary 74.10.** If the absolute Galois group of a field is finite, then it is isomorphic to $1$ or $\mathbb{Z}/2\mathbb{Z}$.

Of course, if $F$ is a field with $\text{char } F = p > 0$ and $\bar{F}$ an algebraic closure, it is possible for the absolute Galois group of $F$ to be trivial, but $\bar{F}/F$ infinite.

**Exercise 74.11.** Prove Facts 74.1.
75. Hilbert’s 17th Problem

In this section, we present a solution to Hilbert’s 17th Problem on positive functions over the reals. We establish the Lang Homomorphism Theorem to solve this problem. We present the proof of the Lang Homomorphism Theorem based on Sylvester’s Lemma 73.10. To use it, we shall need the fact (which we do not prove) that the trace form of an algebra over a field of characteristic zero can be diagonalized, i.e., has a diagonal matrix representation. This is in fact true for all symmetric bilinear forms over fields of characteristic not two. We prove this for the special case that we need.

Lemma 75.1. Let $F$ be a field of characteristic zero and $K/E$ be a finite field extension of degree $n$. Then there exists an $E$-basis $B$ for $K$ such that the matrix representation of the trace form $\varphi : K \times K \to E$ in the basis $B$ is diagonal.

Proof. We must produce an $E$-basis $\{x_1, \ldots, x_n\}$ for $K$ satisfying $\varphi(x_i, x_j) = \delta_{ij}$ for all $i, j = 1, \ldots, n$. As $K$ is a field and $E$ perfect, the trace form is just the non-degenerate field trace $\text{Tr}_{K/E} : K \times K \to E$ by Lemma 77.1. In particular, if $x$ is a nonzero element in $E$, then there exists a $y$ in $E$ with $\text{Tr}_{K/E}(xy)$ nonzero. As $\varphi(1, x) = \delta_{11}$, the restriction of $\text{Tr}_{K/E}$ to $(Ez)^\perp = \{x \mid \text{Tr}_{K/E}(zy) = 0\}$, gives an ‘orthogonal’ decomposition $K = Ez \perp (Ez)^\perp$. Repeating the argument yields the result.

Theorem 75.2. (Lang Homomorphism Theorem) Let $F$ be a real closed field and $K = F(x_1, \ldots, x_n)$ a proper finitely generated formally real field extension of $F$ (hence not algebraic). Then there exists infinitely many $F$-algebra homomorphisms $F[x_1, \ldots, x_n] \to F$, i.e., if $R$ is an affine $F$-algebra that is a domain, then there exist infinitely many $F$-algebra homomorphisms $R \to F$.

Proof. Let $m = \text{tr deg } K/F$. By hypothesis, $m > 0$. We induct on $m$.

Suppose that $m = 1$. We may assume that $x = x_1$ is transcendental over $F$. As $\text{char } F = 0$, we can write $K = F(x)[y]$ for some $y \in K$ and
Furthermore we may assume that \( y \) is integral over \( F[x] \) (cf. Remark 76.14(1)).

**Claim 75.3.** There exist (infinitely many) \( F \)-algebra homomorphisms \( \sigma : F[x, y] \to F \):

Let \( f = f(x, t) = t^r + a_{r-1}(x)t^{r-1} + \cdots + a_0(x) \) with the \( a_i(x) \in F[x] \) be the minimal polynomial of integral \( y \) over \( F[x] \). It suffices to show that there exist (infinitely many) \((a, b) \in F^2\) such that \( f(a, b) = 0 \), since for each such \((a, b)\) the map \( \sigma : F[x, y] \to F \) via \( x \mapsto a \) and \( y \mapsto b \) is well-defined. Let \( P \in \mathcal{X}(K) \) and let \( \hat{K} \) be a real closure of \((K, P)\). So we have \( P \cap K = K^2 \). Let \( \varphi \) be the trace form of the \( F(x)\)-algebra \( K \). Since \( f \) has a root \( y \) in \( \hat{K} \), by Sylvester’s Lemma 73.10, we have \( \text{sgn}_P \varphi > 0 \). As \( r = [K : F(x)] \) and \( \varphi \) can be diagonalized over \( F(x) \), we may assume that there is an \( F(x)\)-basis \( B \) for \( K \) such that

\[
[\varphi]_B = \begin{pmatrix}
h_1(x) & 0 & \cdots & 0 \\
0 & \ddots & \cdots & 0 \\
\vdots & \ddots & \ddots & \vdots \\
h_r(x) & \cdots & 0 & h_1(x)
\end{pmatrix}
\]

is diagonal for some \( h_i(x) \in F(x) \). Moreover, modifying this basis, we may also assume that all the \( h_i(x) \in F[x] \) and are square-free.

As \( F \) is real closed, each of the \( h_i(x) \) can be factored as a product of a constant in \( F \), monic linear polynomials, and monic irreducible quadratic polynomials in \( x \). Each irreducible monic polynomial \( x^2 + ax + b = (x + \frac{a}{2})^2 + (b - \frac{a^2}{4}) \), a sum of two squares in \( F(x) \) hence totally positive as is \( c^2 + ac + b \) for all \( c \in F \) by substitution. Thus each \( h_i \) is a product of a totally positive or totally negative element and finitely many positive factors \( x - a \) relative to \( P \) and finitely many negative factors \( x - b \) relative to \( P \) for some \( a, b \in F \).

Among all the finitely many \( a \)'s occurring in all the \( h_i \), let \( a_0 \) be the maximum relative to the ordering \( F^2 = P \cap F \) and among all the finitely many \( b \)'s occurring in all the \( h_i \), let \( b_0 \) be the minimum relative to the ordering \( P \cap F \). (If there are no \( a \)'s (respectively, \( b \)'s) choose \( a_0 \) (respectively, \( b_0 \)) such that \( a_0 <_P b_0 \) in \( F \).) Thus we have \( x - a_0 \leq_P x - a \) for all such \( a \) and \( x - b \leq_P x - b_0 \) for all such \( b \). As \( x - b_0 <_P x - a_0 \), we must also have \( a_0 <_P b_0 \). Thus there exist infinitely many \( c \in F \) with \( a_0 <_P c <_P b_0 \) (e.g., \( a_0 + \frac{b_0 - a_0}{2^N} \) for \( N > 0 \)). For each such \( c \), we have \( c - a >_P 0 \) and \( x - a >_P 0 \) for all the \( a \)'s and \( c - b <_P 0 \) and
$x - b <_P 0$ for all the $b$'s. It follows that

$$\text{sgn}_P \varphi = \text{sgn}_{P \cap F} [\varphi]_B = \text{sgn}_{P \cap F} \left( \begin{array}{ccc} h_1(c) & 0 \\ 0 & \ddots \\ & & h_r(c) \end{array} \right)$$

**Assertion.** Let $c \in F$ satisfy $a_0 <_P c <_P b_0$ and $\varphi_c$ be the trace form of the $F$-algebra $A = F[t]/(f(c,t))$. Then, except for finitely many exceptions, there exists an $F$-basis $C$ for $A$ satisfying

$$(*) \quad [\varphi_c]_c = \left( \begin{array}{ccc} h_1(c) & 0 \\ 0 & \ddots \\ & & h_r(c) \end{array} \right).$$

In particular, for each such $c$, we have $\text{sgn}_{P \cap F} \varphi_c > 0$.

As $f(c,t)$ has a root in $K$, hence $\tilde{K}$, if $(*)$ holds then $\text{sgn}_{P \cap F} \varphi_c > 0$ by Sylvester’s lemma 73.10. So we need only show the first part.

Let $B = (b_{ij}(x))$ be the matrix representation of the trace form $\varphi$ relative to the basis $\{1, y, \ldots, y^{r-1}\}$ of $K$ over $F(x)$. As $y$ is integral over $F(x)$, each $b_{ij}(x) \in F[x]$. There exists a change of basis matrix $C = (c_{ij}(x)) \in \text{GL}_r(K)$ such that

$$C^t B C = \left( \begin{array}{ccc} h_1(x) & 0 \\ 0 & \ddots \\ & & h_r(x) \end{array} \right).$$

If $c$ is not a zero of the denominator of any of the $c_{ij}(x)$'s or a zero of $\det C$, then

$$\left. C^t B \right|_{x=c} = \left. C \right|_{x=c} = \left( \begin{array}{ccc} h_1(c) & 0 \\ 0 & \ddots \\ & & h_r(c) \end{array} \right)$$

as desired. As there exist only finitely many $c$ not satisfying these conditions, the assertion follows.

As $F$ is real closed, for each $c$ in the assertion satisfying $\text{sgn}_{P \cap F} \varphi_c > 0$, there exists a $b \in F$ such that $f(c,b) = 0$ by Sylvester’s Lemma 73.10. Thus there exist infinitely many $(c,b)$ such that $f(c,b) = 0$ and the Claim is established.

Next we show that almost all of the $F$-algebra homomorphisms $\sigma : F[x,y] \to F$ extend to $F[x,y,x_2,\ldots,x_n] \to F$. As the transcendence degree $m = \text{tr} \deg_F K = 1$ and $F(x)[y] = F(x_1,\ldots,x_n)$, we can write

$$x_i = \frac{g_i(x,y)}{q_i(x)}$$

for some $g_i \in F[x,t]$, $0 \neq q_i \in F[x]$. 

Let \( q = \prod_i q_i \). As \( q(c) \neq 0 \) for almost all of the \( c \) constructed above, \( \sigma(q(x)) \neq 0 \) for almost all the \( \sigma \) constructed above. Let \( \sigma \) be any such one. By properties of localization with \( S = \{ t^i \mid i \geq 0 \} \), any such \( \sigma \) extends to an \( F\)-algebra homomorphism \( \sigma : F[x,y][\frac{1}{q}] \to F \) by Exercise 28.3(2). Since \( F[x_1,\ldots,x_n] \subset F[x_1,y,x_2,\ldots,x_n] \subset F[x,y][\frac{1}{q}] \), the \( m = 1 \) case is completed.

Suppose that \( m > 1 \). Choose \( K/E/F \) satisfying \( \text{tr deg } E/K = 1 \). Let \( \tilde{E} \) be the field of algebraic elements over \( E \) lying in a real closure \( \tilde{K} \) of \( K \) (called the algebraic closure of \( E \) in \( \tilde{K} \)). By the Fundamental Theorem of Algebra, \( \tilde{E} \) is also real closed. As \( \text{tr deg } \tilde{E}(x_1,\ldots,x_n) = 1 \), there exist infinitely many \( \tilde{E} \)-algebra homomorphisms \( \sigma : \tilde{E}(x_1,\ldots,x_n) \to \tilde{E} \) by the \( m = 1 \) case. Let \( \sigma \) be any one of these. If \( \sigma(x_i) \in F \) for all \( i \), then \( \sigma|_{F[x_1,\ldots,x_n]} : F[x_1,\ldots,x_n] \to F \) is an \( F \)-algebra homomorphism.

So suppose that there exists an \( i \) with \( \sigma(x_i) \not\in F \). As \( \sigma(x_i) \in \tilde{E} \) for all \( i \), we have \( \text{tr deg } F(\sigma(x_1),\ldots,\sigma(x_n)) \leq \text{tr deg } \tilde{E} = \text{tr deg } E < m \). By induction there exist infinitely many \( F \)-algebra homomorphisms \( \tau : F[\sigma(x_1),\ldots,\sigma(x_n)] \to F \). Then the compositions \( \tau \circ \sigma|_{F[x_1,\ldots,x_n]} : F[x_1,\ldots,x_n] \to F \) give infinitely many \( F \)-algebra homomorphisms. \( \square \)

A commutative ring is called semi-real if \(-1 \) is not a sum of squares in \( R \) and is called formally real if \( \sum x_i^2 = 0 \) in \( R \) implies that \( x_i = 0 \) for all \( i \). More generally, an ideal \( \mathfrak{A} \) in a commutative ring is called semi-real (respectively, formally real) if \( R/\mathfrak{A} \) is semi-real (respectively, formally real). We note that if \( R \) is a domain, then it is formally real if and only if its quotient field is and a field is semi-real if and only if formally real.

**Lemma 75.4.** Let \( R \) be a commutative ring. Then the following are equivalent

1. \( R \) is semi-real.
2. There exists \( P \in \mathcal{Y}(R) \) such that \( P \cap -P \) is a prime ideal.
3. There exists a formally real prime ideal \( \mathfrak{p} \) in \( R \).
4. There exists a ring homomorphism \( R \to F \) for some formally real field \( F \).

**Proof.** (1) \( \Rightarrow \) (2): As \( R \) is semi-real, we have \( \sum R^2 \in \mathcal{Y}(R) \). Thus (2) follows by Lemma 72.4.

(2) \( \Rightarrow \) (3): The prime ideal \( P \cap -P \) in (2) is clearly semi-real and excludes the multiplicative set \( S = 1 + \sum R^2 \). By choice \( 0 \not\in S \) and \( S + \sum R^2 \subset S \). By Zorn’s lemma there exists an ideal \( \mathfrak{p} \) containing
Corollary 75.5. Let $F$ be a real closed field and $R$ a domain that is an affine $F$-algebra. If $R$ is semi-real the there exists an $F$-algebra homomorphism $R \to F$.

**Proof.** There exists a formally real prime ideal $p$ in $R$. Then $R/p$ is a real $F$-affine ring that is a domain so there exists an $F$-algebra homomorphism $R/p \to R$ by the Lang Homomorphism Theorem 75.2. The composition $R \to R/p \to R$ now works. □

Corollary 75.6. Let $F$ be a real closed field and $R$ a real affine $F$-algebra that is a domain. Let $f_1, \ldots, f_n \in R \setminus \{0\}$. Then there exists an $F$-algebra homomorphism $\varphi : R \to F$ so that $\varphi(f_i) \neq 0$ for all $i$.

**Proof.** $R[f_1^{-1}, \ldots, f_n^{-1}]$ is a real affine $F$-algebra that is a domain. □

Note that $R = \mathbb{R}[t_1, \ldots, t_n]/(t_1^2 + \cdots + t_n^2)$ is a semi-real affine $\mathbb{R}$-algebra that is a domain but it is not formally real since any $\mathbb{R}$-algebra homomorphism $R \to \mathbb{R}$ must take each $t_i \mapsto 0$.

Corollary 75.7. Let $F$ be a real closed field and $R$ an affine $F$-algebra. Let $f_1, \ldots, f_r, g_1, \ldots, g_s$ lie in $R$. Suppose that there exists a maximal preordering $P \in \mathcal{Y}(R)$ such that $0 \neq f_i \in P$ for all $i$ and $g_j \in P$ for all $j$. Then there exists an $F$-algebra homomorphism $\varphi : R \to F$ so that $0 \neq \varphi(f_i) \in F^2$ and $\varphi(g_j) \in F^2$ for all $i$ and $j$.

**Proof.** Let $p = P \cap -P$, a prime ideal in $R$. Replacing $R$ by $R/p$ we may assume that $P \cap -P = p = 0$. In particular, $R$ is a domain and $P$ is an ordering on $R$. Clearly, $P$ extends to an ordering of the quotient field $K$ of $R$. Also call this extension $P$. We have $f_i \geq_P 0$ and $g_j \geq_P 0$ for all $i$ and $j$. Let $E = K(\sqrt{f_1}, \ldots, \sqrt{f_r}, \sqrt{g_1}, \ldots, \sqrt{g_s})$. By induction on $r + s$, the ordering $P$ extends to an ordering $Q$ on $E$. Since the domain $A = R[f_1^{-1}, \ldots, f_r^{-1}, \sqrt{f_1}, \ldots, \sqrt{f_r}, \sqrt{g_1}, \ldots, \sqrt{g_s}]$ is a real affine $F$-algebra, there exists an $F$-algebra homomorphism $A \to F$. This homomorphism has the desired properties. □

If $\mathfrak{A}$ be an ideal in $F[t_1, \ldots, t_n]$ with $F$ infinite, let

1. $Z_F(\mathfrak{A}) = \{ x \in F^n \mid f(x) = 0 \text{ for all } f \in \mathfrak{A} \}$.
2. $A_F(\mathfrak{A}) = F[t_1, \ldots, t_n]/\mathfrak{A}$. 
Suppose that $p = \mathfrak{a}$ is a prime ideal in $F[t_1, \ldots, t_n]$. So $A_F(p)$ is a domain. As usual we say a rational function $f$ in the quotient field of $A_F(p)$ is defined at $x \in Z_F(p)$ if $f = g/h$ for some $g, h \in A_F(\mathfrak{a})$ with $h(x) \neq 0$. If $F$ is real closed, a rational function $f$ in the quotient field of $A_F(p)$ is called positive semi-definite if $f(x) \geq 0$ for all points $x \in A_F(p)$ where $f$ is defined.

**Corollary 75.8.** (Weak Real Nullstellensatz) If $F$ is real closed and $\mathfrak{a}$ is an ideal in $F[t_1, \ldots, t_n]$, then $\mathfrak{a}$ is semi-real if and only if $Z_F(\mathfrak{a}) \neq \emptyset$.

**Proof.** Let $A_F(\mathfrak{a}) = F[x_1, \ldots, x_n]$. If $(a_1, \ldots, a_n) \in Z_F(\mathfrak{a})$, then $f : A_F(\mathfrak{a}) \to F$ by $x_i \mapsto a_i$ defines an $F$-algebra homomorphism and $\mathfrak{a}$ is semi-real by Lemma 75.4. Suppose that $\mathfrak{a}$ is semi-real. By Lemma 75.4 there exists a prime ideal $p$ in $F[t_1, \ldots, t_n]$ such that $p/\mathfrak{a}$ is a formally real prime ideal in $A_F(\mathfrak{a})$. It follows that $F[t_1, \ldots, t_n]/p$ is a formally real domain so by the Lang Homomorphism Theorem 75.2 there exists an $F$-algebra homomorphism $\sigma : F[t_1, \ldots, t_n]/p \to F$. Then $(a_1, \ldots, a_n) \in Z_F(\mathfrak{a})$ with $a_i = \sigma(t_i + p)$. \hfill \square

**Theorem 75.9.** (Artin) Let $F$ be a real closed field and $p$ a formally real prime ideal in $F[t_1, \ldots, t_n]$. Let $K$ be the quotient field of $A_F(p)$. Let $f \in K$. If $f$ is positive, then $f$ is a sum of squares in $K$.

**Proof.** Suppose that $f$ is not a sum of squares in $K$. Then there exists $P \in \mathcal{X}(K)$ such that $f <_P 0$ by Corollary 72.6. Thus $P$ extends to $E = K(\sqrt{-f})$ by Theorem 73.2. Let $f = g/h$, with $g, h \in A_F(p)$ and let $- : F[t_1, \ldots, t_n] \to A_F(p)$ be the canonical map. By the Lang Homomorphism Theorem, there exists an $F$-algebra homomorphism $\sigma : F[t_1, \ldots, t_n, \sqrt{-g}, \frac{1}{gh}] \to F$. For each $i$, let $a_i = \sigma(\overline{t}_i)$. If $q \in p$ then

$$q(a_1, \ldots, a_n) = q(\sigma(\overline{t}_1), \ldots, \sigma(\overline{t}_n)) = \sigma(q(\overline{t}_1, \ldots, \overline{t}_n)) = 0,$$

so $(a_1, \ldots, a_n) \in Z_F(p)$. As $\sigma(g) \neq 0$ and $\sigma(h) \neq 0$, the rational function $f$ is defined at $(a_1, \ldots, a_n)$ and is not zero at it. Since $\sigma(-f) \in F^2$, we must have $\sigma(f) <_P 0$. But $\sigma(f) = f(a_1, \ldots, a_n) >_P 0$, a contradiction. \hfill \square

Artin’s Theorem immediately answers Hilbert’s 17th problem whether every positive semi-definite function $f$ in $\mathbb{R}(t_1, \ldots, t_n)$ is a sum of squares in $K$ in the affirmative.

**Corollary 75.10.** (Hilbert’s 17th Problem) Any positive semi-definite function in $\mathbb{R}(t_1, \ldots, t_n)$ is a sum of squares in $\mathbb{R}(t_1, \ldots, t_n)$.

The Motzkin polynomial $t_1^4t_2^2 + t_1^2t_2^4 - 3t_1^2t_2^2 + 1$ in $\mathbb{R}[t_1, t_2]$ is positive semi-definite but not a sum of squares in $\mathbb{R}[t_1, t_2]$. [It is a sum
of four squares, but not three squares, in $\mathbb{R}(t_1, t_2)$. Cf. Lam’s book *Introduction to quadratic forms over fields* [33] for more detail.]
CHAPTER XV

Dedekind Domains

The study of fields arose from the desire to understand the roots of polynomials with rational coefficients. In a similar way, algebraic number theory arose from the desire to understand the roots of monic polynomials with integer coefficients. In this chapter, we give an introduction to algebraic number theory. Since the generalization to studying Dedekind domains – an appropriate generalization of the integers – is not too difficult, we take this more general approach.

76. Integral Elements

We begin with the study of monic polynomials over (nonzero) commutative rings.

**Definition 76.1.** Let $B$ be a nonzero commutative ring and $A$ a subring of $B$. We shall write this as $B/A$ and call it a *ring extension* mimicking how we wrote field extensions. An element $x$ in $B$ is called *integral over* $A$ if $x$ is a root of a monic polynomial $f$ in $A[t]$.

**Examples 76.2.**
1. Let $K/F$ be an extension of fields. Then an element $x$ in $K$ is integral over $F$ if and only if it is algebraic over $F$.
2. If $A$ is a nonzero commutative ring, then every element $x$ in $A$ is integral over $A$ as it is a root of $t - x$ in $A[t]$.
3. Let $d$ be a square-free integer and $x = a + b\sqrt{d}$ with $a$ and $b$ in $\mathbb{Q}$. Then $x$ is a root of $t^2 - 2at + (a^2 - b^2d)$. In particular, one checks that $x$ is integral over $\mathbb{Z}$ if and only if $2a$ and $a^2 - b^2d$ are integers. For example, the complex numbers $\sqrt{2}$ and $(-1 + \sqrt{-3})/2$ are integral over $\mathbb{Z}$.
4. Every $n$th root of unity in $\mathbb{C}$, $n > 0$, is integral over $\mathbb{Z}$.
5. Let $C/B/A$ be ring extensions. If $x$ in $C$ is integral over $A$, then it is integral over $B$.

Because of the first example, we expect that some of our field theoretic ideas and results should generalize. This is true. First, we must
find a ring theoretic test for an element to be integral. It is given by
the following basic result.

**Proposition 76.3.** Let $B/A$ be an extension of nonzero commutative
rings and $x$ an element in $B$. Then $x$ is integral over $A$ if and only if
$A[x]$ is a finitely generated $A$-module.

**Proof.** ($\Rightarrow$): Let $x$ be a root of the monic polynomial $f$ in $A[t]$.
By the General Division Algorithm 33.4, we can write
$$g = fq + r$$
with $q, r \in A[t]$ and $r = 0$ or $\deg r < \deg f$.
Therefore, if $n = \deg f$, we have
$$g(x) = r(x) \in \sum_{i=0}^{n-1} Ax^i,$$
so $A[x] = \sum_{i=0}^{n-1} Ax^i$.

($\Leftarrow$): We prove a more general statement.

**Claim 76.4.** Let $M$ be an $A[x]$-module that is finitely generated as an
$A$-module. If $\text{ann}_{A[x]} M = 0$ (i.e., if $r$ in $A[x]$ satisfies $rm = 0$ for all $m$
in $M$, then $r = 0$), then $x$ is integral over $A$:
Suppose that $M = Am_1 + \cdots + Am_n$. By assumption, $M$ is an $A[x]$-
module, i.e.,
$$\lambda_x : M \to M \text{ given by } m \mapsto xm$$
is an $A$-endomorphism. For each $i$, $i = 1, \ldots, n$, we can write
$$xm_i = \sum_{j=1}^{n} a_{ij}m_j \text{ some } a_{ij} \in A, 1 \leq i, j \leq n.$$ Therefore,
$$\sum_{j=1}^{n} (\delta_{ij}x - a_{ij})m_j = 0 \text{ for } i = 1, \ldots, n.$$ Let $\Delta$ be the determinant of the $n \times n$ matrix $(\delta_{ij}x - a_{ij})$. By Cramer’s
Rule, $\Delta m_i = 0$ for $i = 1, \ldots, n$, so $\Delta M = 0$, i.e., $\Delta$ lies in $\text{ann}_{A[x]} M = 0$. It follows that $x$ is a root of the monic polynomial $\det((\delta_{ij}t - a_{ij}))$
in $A[t]$. This proves the claim.

Let $M = A[x]$, a finitely generated $A$-module by assumption. As $A[x]$ is
a subring of $B$ it has a one (not zero). Consequently, $\text{ann}_{A[x]} A[x] = 0$,
so the claim yields the desired result.

**Corollary 76.5.** Let $B/A$ be an extension of nonzero commutative
rings, $x_1, \ldots, x_n$ elements of $B$. Then $x_1, \ldots, x_n$ are all integral over
$A$ if and only if $A[x_1, \ldots, x_n]$ is a finitely generated $A$-module.
Proof. \((\Rightarrow)\): The case for \(n = 1\) follows from the proposition. By induction, \(A[x_1, \ldots, x_{n-1}]\) is a finitely generated \(A\)-module, say with generating set \(S\). As \(x_n\) is integral over \(A\), it is a root of a monic polynomial in \(A[t]\), say of degree \(d\). Then \(\bigcup_{i=0}^{d-1}\{sx^i \mid s \in S\}\) is a generating set for \(A[x_1, \ldots, x_n]\).

\((\Leftarrow)\): Since \(\text{ann}_{A[x_i]} A[x_1, \ldots, x_n] = 0\) for \(i = 1, \ldots, n\), this follows from Claim 76.4. \(\square\)

In the above, we could also have used the following lemma that we leave as an exercise. We shall need it below.

**Lemma 76.6.** Let \(B/A\) be an extension of nonzero commutative rings. If \(M\) is a finitely generated \(B\)-module and \(B\) a finitely generated \(A\)-module, then \(M\) is a finitely generated \(A\)-module.

Note the corollary above is analogous with the result that when \(K/F\) is a field extension, \(x_1, \ldots, x_n\) elements in \(K\), then \(F(x_1, \ldots, x_n)/F\) is finite if and only if \(x_1, \ldots, x_n\) are all algebraic over \(F\). In a similar way we have the following analogue of another field theoretical result.

**Corollary 76.7.** Let \(B/A\) be an extension of nonzero commutative rings, \(x\) and \(y\) elements of \(B\) both integral over \(A\). Then \(x \pm y\), and \(xy\) are integral over \(A\). In particular,

\[\{x \in B \mid x \text{ integral over } A\}\]

is a subring of \(B\).

Proof. If \(z\) is any element in \(A[x, y]\), then \(A[x, y, z] = A[x, y]\) is a finitely generated \(A\)-module, so \(z\) is integral over \(A\). \(\square\)

Of course, if \(B/A\) is an extension of commutative rings with \(x\) in \(B\) a unit in \(B\), it may be integral over \(A\) but its inverse not. For example, \(1/2\) is not integral over \(\mathbb{Z}\).

**Definition 76.8.** Let \(B/A\) be an extension of nonzero commutative rings and

\[C = \{x \in B \mid x \text{ integral over } A\}\].

The ring \(C\) is called the integral closure of \(A\) in \(B\). If \(C = B\), i.e., every element of \(B\) is integral over \(A\), we say \(B/A\) is integral and if \(C = A\), we say that \(A\) is integrally closed in \(B\).

Therefore, \(B/A\) being integral is the analogue of an algebraic extension in field theory. Analogous to the field case, we have:

**Corollary 76.9.** Let \(C/B\) and \(B/A\) be extensions of nonzero commutative rings. Then \(C/B\) and \(B/A\) are integral extensions if and only if \(C/A\) is an integral extension.
Proof. We need only show if $C/B$ and $B/A$ are integral extensions so is $C/A$. Let $c$ be an element in $C$. As $c$ is integral over $B$, it satisfies an equation $c^n + b_{n-1}c^{n-1} + \cdots + b_0 = 0$ in $B$ for some $b_0, \ldots, b_{n-1}$ in $B$. Let $B_0 = A[b_0, \ldots, b_{n-1}]$. As $B/A$ is integral, it follows that $A[c, b_0, \ldots, b_{n-1}]$ is a finitely generated $B_0$-module by Proposition 76.3 and $B_0$ is a finitely generated $A$-module by Corollary 76.5. The result now follows using Lemma 76.6. □

Corollary 76.10. Let $C/B/A$ be extensions of nonzero commutative rings with $B/A$ integral. Then the integral closure of $A$ in $C$ is the same as the integral closure of $B$ in $C$.

We shall mostly be interested in the case that our rings are domains. For this case, we shall use the following notation:

Notation 76.11. If $A$ is a domain and $K$ any field containing $A$, we shall denote the integral closure of $A$ in $K$ by $A_K$, i.e.,

$$A_K := \{ x \in K \mid x \text{ is integral over } A \}.$$ 

Definition 76.12. Let $A$ be a domain. We say that $A$ is integrally closed or a normal domain if $A$ is integrally closed in its quotient field, i.e., if $F$ is the quotient field of $A$, then $A = A_F$.

An important example of integrally closed domains is given by:

Proposition 76.13. A UFD is integrally closed. In particular, any PID is integrally closed.

Proof. Let $A$ be a UFD and $a$ and $b$ be elements in $A$ with $b$ nonzero. Suppose that $a/b$ is integral over $A$. As $A$ is a UFD, we may assume that $a$ and $b$ are relatively prime. By definition, $a/b$ satisfies an equation

$$(\frac{a}{b})^n + c_1(\frac{a}{b})^{n-1} + \cdots + c_n = 0$$

for some $c_1, \ldots, c_n$ in $A$. Therefore, we have the equation

$$a^n = -b(c_1a^{n-1} + \cdots + b^{n-1}c_n) \text{ in } A.$$ 

It follows that $b | a^n$ in the UFD $A$, hence $b$ is a unit in $A$ and $a/b$ lies in $A$. □

Remarks 76.14. Let $A$ be a domain, $F$ its quotient field, and $L/K/F$ algebraic field extensions.

1. If $x$ lies in $L$, then there exists a nonzero element $c$ in $A$ satisfying $cx$ is integral over $A$, i.e., lies in $A_L$. Indeed if $x$ satisfies

$$a_nx^n + a_{n-1}x^{n-1} + \cdots + a_0 = 0$$
in $L$ with $a_0, \ldots, a_n$ elements of $A$ and $a_n$ nonzero, then multiplying this equation by $a_n^{n-1}$ shows that $a_n x$ is integral over $A$.

2. We have $A_L \cap K = A_K$.

3. The quotient field of $A_K$ is $K$ by the first remark.

4. If $E/F$ is a field extension and $E/K/F$ with $K$ the maximal algebraic field extension of $F$ in $E$, then $A_K = A_L$.

The case in which we are most interested is when $A$ is the ring of integers. If $K/\mathbb{Q}$ is a finite field extension, then $K$ is called an algebraic number field and $\mathbb{Z}_K$ is called the ring of algebraic integers in $K$. Let $\overline{\mathbb{Q}}$ be the algebraic closure of $\mathbb{Q}$ in $\mathbb{C}$. Then algebraic number theory is the study of $\mathbb{Z}_{\overline{\mathbb{Q}}}$. As $\overline{\mathbb{Q}}$ is the union of number fields, we have $\mathbb{Z}_K = \mathbb{Z}_{\overline{\mathbb{Q}}} \cap K$ for any number field $K$. Since $\mathbb{Q}$ is perfect, we shall not have to worry about separability that in the general case can cause serious problems.

**Exercises 76.15.**

1. Let $d$ be a square-free integer and $x = a + b \sqrt{d}$ with $a$ and $b$ in $\mathbb{Q}$. Show that $x$ is integral over $\mathbb{Z}$ if and only if $2a$ and $a^2 - b^2 d$ are integers.

2. Let $\varphi : A \to B$ be a ring homomorphism and $M$ a finitely generated $B$-module. If $B$ is a finitely generated $A$-module via the pullback, then $M$ is a finitely generated $A$-module. In particular, Lemma 76.6 is valid.

3. Let $\varphi : A \to B$ be a ring homomorphism of commutative rings. We say that $\varphi$ is integral if $B/\varphi(A)$ is integral. Show that a composition of integral homomorphisms is integral.

4. Let $\varphi : A \to B$ be a ring homomorphism of commutative rings. We say that $\varphi$ is finite if $B$ is a finitely generated $A$-module (via $\varphi$) and of finite type if there exists a ring epimorphism $A[t_1, \ldots, t_n] \to B$ for some $n$ lifting $\varphi$. Show that $\varphi$ is finite if and only if $\varphi$ is integral and of finite type.

5. Let $A$ be a domain with quotient field $F$ and $K/F$ a field extension. If $S$ is a multiplicative set (not containing zero) in $A$, show that $(S^{-1}A)_K = S^{-1}(A_K)$, i.e., taking localization and integral closures commute. In particular, this applies when $S = A \setminus \mathfrak{p}$ with $\mathfrak{p}$ a prime ideal in $A$.

**77. Integral Extensions of Domains**

In this section, we shall investigate the integral closure of a domain in a separable extension of its quotient field. Separability makes
life much nicer, and since we are mostly interested in applications to algebraic number theory, this will not be an impairment.

When we investigated cyclic extensions of fields, we needed to use the norm of such an extension. Although we introduced the notion of trace at the same time, except for exercises, we did not use it. Although the norm and trace of an arbitrary, i.e., not necessarily separable, finite extension can be defined, this is not useful for the trace. Indeed, the trace from a finite inseparable extension turns out to be the zero map. However, for the separable case, the trace is very useful. This was because of Exercise 59.26(5). Because this is important to our current investigation, we now prove this exercise.

Lemma 77.1. Let $K/F$ be a finite separable extension of fields of degree $n$, $K^* := \text{Hom}_F(K,F)$, the $F$-linear dual space of $K$. Then the map $T : K \to K^*$ defined by $x \mapsto \text{tr}_x : y \mapsto \text{Tr}_{K/F}(xy)$ is an $F$-linear isomorphism. In particular, if $B = \{w_1, \ldots, w_n\}$ is an $F$-basis for $K$, then there exists an $F$-basis $\{w'_1, \ldots, w'_n\}$ for $K$ satisfying $\text{Tr}_{K/F}(w_iw'_j) = \delta_{ij}$ for all $i,j = 1, \ldots, n$.

Remark 77.2. The basis $\{w'_1, \ldots, w'_n\}$ in the lemma is called the complementary basis to $B$.

Proof. By Dedekind’s Lemma 53.3, the map $\text{Tr}_{K/F} : K \to F$ is nontrivial, so there exists an element $z$ in $K$ with $\text{Tr}_{K/F}(z)$ nonzero. Therefore, $\text{Tr}_{K/F}(xx^{-1}z)$ is not zero for every nonzero $x$ in $K$. It follows that $\text{tr}_x$ is nonzero for all nonzero $x$ in $K$, i.e., $T : K \to K^*$ is injective, hence an isomorphism by dimension count. Let $\{f_1, \ldots, f_n\}$ be the dual basis to $B$ and choose $w'_i$ in $K$ to satisfy $T(w'_i) = f_i$ for $i = 1, \ldots, n$. Then $\{w'_1, \ldots, w'_n\}$ works. □

Suppose that $K_i/F$ is a field extension for $i = 1, 2$ and $\sigma : K_1 \to K_2$ is an $F$-isomorphism. Let $A$ be a domain lying in $F$. If $x$ in $K_1$ is integral over $A$ then $\sigma(x)$ is also integral over $A$, i.e., $\sigma(A_{K_1}) \subset A_{K_2}$. In particular, if $K = K_1 = K_2$ then $\sigma(A_K) = A_K$. This means that we can use Galois theory. We shall need the results about the trace (and norm) for finite separable extensions as stated in Remark 59.18 as well as Exercise 59.26(3) which we now prove.

Lemma 77.3. Let $K/F$ be a finite separable extension of fields. Then $N_{K/F} = N_{E/F} \circ N_{K/E}$ and $\text{Tr}_{K/F} = \text{Tr}_{E/F} \circ \text{Tr}_{K/E}$ for any intermediate field $K/E/F$.

Proof. Let $L/F$ be a finite Galois extension with $L/K$. Let $\sigma_1, \ldots, \sigma_n : E \to L$ denote all the $F$-homomorphisms and $\tau_1, \ldots, \tau_m : K \to L$ all the $E$-homomorphisms. Extend each $\sigma_i$ to $\tilde{\sigma}_i$ in $G(L/F)$.
If $\mu : K \to L$ is an $F$-homomorphism, then $\mu|_E = \sigma_i$ for some $i$ and $\sigma_i^{-1}\mu = \tau_j$ for some $j$, hence $\mu = \sigma_i\tau_j$. As in the proof of Proposition 55.11, we see that these are all the distinct $E$-homomorphisms $K \to L$. It follows that if $x$ lies in $K$, then
\[
\text{Tr}_{K/F}(x) = \sum \sigma_i \tau_j(x) = \sum \sigma_i(\text{Tr}_{K/E}(x)) = \text{Tr}_{E/F}(\text{Tr}_{K/E}(x))
\]
and similarly for the norm. \hfill \Box

Let $K/F$ be a finite separable extension with $L/K$ a normal closure of $K/F$ and $\sigma_1, \ldots, \sigma_n : K \to L$ all the $F$-homomorphisms. If $x$ is an element of $K$, then any elementary symmetric function in $\sigma_1(x), \ldots, \sigma_n(x)$ is fixed by $G(L/F)$ hence lies in $F$. It also follows that $N_{K/F}(x) = x[\sigma:K:F(x)]$.

**Proposition 77.4.** Let $A$ be a domain with quotient field $F$ and $K/F$ a finite separable extension. If $x$ in $K$ is integral over $A$, then the minimal polynomial $m_F(x)$ of $x$ lies in $A_F[t]$ and both $\text{Tr}_{K/F}(x)$ and $N_{K/F}(x)$ lie in $A_F$. In particular, if $A$ is integrally closed, then $m_F(x)$ lies in $A[t]$ and $\text{Tr}_{K/F}(x)$ and $N_{K/F}(x)$ lie in $A$.

**Proof.** Let $L/F$ be a finite Galois extension satisfying $L/K$ and $\sigma_1, \ldots, \sigma_n : K \to L$ all the distinct $F$-homomorphisms. If $x$ is an element of $K$, then any elementary symmetric function in $\sigma_1(x), \ldots, \sigma_n(x)$ is fixed by $G(L/F)$ hence lies in $F$. Therefore, $m_F(x) = \prod_{i=1}^n (t - \sigma_i(x))$ lies in $A_L[t] \cap F[t]$. It follows that $m_F(x)$ lies in $A_F[t]$. By the lemma and Property 59.17(3) (cf. Remark 59.18), we have
\[
N_{F(x)/F}(N_{K/F}(x)) = (N_{F(x)/F}(x))^{[K:F(x)]},
\]
hence lies in $A_L \cap F = A_F$, and similarly for the trace. \hfill \Box

The result for which we are aiming can now be established.

**Theorem 77.5.** Let $A$ be an integrally closed Noetherian domain with quotient field $F$. Let $K/F$ be a finite separable extension. Then the integral closure of $A$ in $K$ is a finitely generated $A$-module. In particular, $A_K$ is also an integrally closed Noetherian domain.

**Proof.** By Theorem 39.7, we know that any finitely generated $A$-module is a Noetherian $A$-module, since $A$ is a Noetherian ring. In particular, it suffices to show there exists a finitely generated $A$-module $M$ with $A_K$ a submodule of $M$. Let $n = [K : F]$ and $B = \{w_1, \ldots, w_n\}$ be an $F$-basis for $K$. As $K$ is the quotient field of $A_K$, by clearing denominators, we may assume that $B$ lies in $A_K$. As $K/F$ is a finite separable extension, there exists a complementary basis $B' = \{w'_1, \ldots, w'_n\}$ to $B$ for $K$ by Lemma 77.1, i.e., $\text{Tr}_{K/F}(w_iw'_j) = \delta_{ij}$ for all
i and j. For each j, j = 1, . . . , n, there exists a nonzero c_j in A such that c_j w'_j lies in A_K by Remark 76.14(1). Let c = c_1 · · · c_n ≠ 0. To finish, it suffices to show the following:

Claim 77.6. A_K is a submodule of \( \sum A c_w w_i \):

Let z be an element of A_K. We can write 
\[ z = \sum_{i=1}^{n} b_i w_i \]
for some b_1, . . . , b_n in F. Then for each j, 1 ≤ j ≤ n, we have
\[ \text{Tr}_{K/F}(czw'_j) = c \text{Tr}_{K/F}(zw'_j) = cb_j. \]
As z and each cw'_j lies in A_K, we conclude that cb_j lies in A_F = A for every j, i.e., b_j ∈ c^{-1}A for each j. The result follows. □

Remark 77.7. The result is false if we drop the separability assumption. Indeed there exist counterexamples with A a PID with precisely one nonzero prime ideal.

Corollary 77.8. Let A be a PID with quotient field F and K/F a finite separable field extension. Then A_K is a finitely generated free A-module of rank \([K : F]\).

PROOF. By the theorem, A_K is a finitely generated A-module. Since A_K is a domain, it is a torsion-free A-module, hence A-free by Corollary 43.16 to the Fundamental Theorem of Finite Generated Modules over a PID. Therefore, we need only compute the rank of A_K. We showed in the proof above that there exists an F-basis for K lying in A_K, so the rank of A_K is at most \([K : F]\) using Proposition 43.1. But any A-linearly independent set in A_K is F-linearly independent, so we must have A_K be of rank \([K : F]\). □

Let A be a domain with quotient field F and K/F a finite field extension. If A_K is a free A-module of rank \([K : F]\), then an A-basis for is called an integral basis. The corollary applies to any ring of algebraic numbers \( \mathbb{Z}_K \). It says that it is a free abelian group of rank \([K : \mathbb{Q}]\), hence has an integral basis. The theorem also applies to the case that A = F[t] with F a field, and K/F(t) a finite separable extension. Then F[t]_K is a free F[t]-module of rank \([K : F(t)]\). These are the primary examples of our next study, when we further restrict our domain.

Exercise 77.9. Let K be a quadratic extension of \( \mathbb{Q} \). Write K = \( \mathbb{Q}(\sqrt{D}) \) with D a square-free integer. Show every element \( \alpha \) of K can be written as \( \alpha = (a + b\sqrt{D})/2 \) with a and b integers. In particular, \( \text{Tr}_{K/\mathbb{Q}}(\alpha) = a \) and \( N_{K/\mathbb{Q}}(\alpha) = (a^2 + b^2 D)/4 \). Show
78. Dedekind Domains

We know that PIDs are UFDs, but rings of algebraic integers are not always UFDs, e.g., $\mathbb{Z}[\sqrt{-5}]$. We wish to study the appropriate generalization of PIDs that allows us to replace the property of UFDs, which is a property about elements, with an analogous property about ideals. We shall begin using the theory that we have already developed. We need the following observation:

**Lemma 78.1.** Let $\varphi : A \to B$ be a ring homomorphism of commutative rings. If $\mathfrak{P}$ is a prime ideal in $B$, then $\varphi^{-1}(\mathfrak{P})$ is a prime ideal in $A$. In particular, if $\varphi$ is the inclusion of rings, then $\mathfrak{P} \cap A$ is a prime ideal in $A$.

**Proof.** The homomorphism $\varphi$ induces a monomorphism of rings $\overline{\varphi} : A/\varphi^{-1}(\mathfrak{P}) \to B/\mathfrak{P}$. As $B/\mathfrak{P}$ is a domain so is $A/\varphi^{-1}(\mathfrak{P})$ and the result follows. \hfill \Box

We come to the main subject of this chapter.

**Definition 78.2.** Let $A$ be a domain, not a field. Then $A$ is called a **Dedekind domain** if $A$ is a Noetherian integrally closed domain in which every nonzero prime ideal is a maximal ideal. Let $\text{Max}(A) = \{ \mathfrak{p} \mid \mathfrak{p} \text{ a maximal ideal} \}$. Then $\text{Spec}(A) = \{0\} \cup \text{Max}(R)$ in a Dedekind domain $A$.

**Examples 78.3.** The verifications of the following are left as exercises:

1. Let $A$ be a PID that is not a field. Then $A$ is a Dedekind domain. In particular, $\mathbb{Z}$ and $\mathbb{F}[t]$, with $\mathbb{F}$ a field, are Dedekind domains.

2. Let $A$ be a Dedekind domain and $S$ a multiplicative set not containing zero. Then the localization $S^{-1}A$ of $A$ is a Dedekind domain or a field. In particular, if $\mathfrak{p}$ is a maximal ideal in $A$, then the localization $A_\mathfrak{p} = S^{-1}A$ with $S = A \setminus \mathfrak{p}$ is a Dedekind domain.

3. Recall that a **local ring** is a commutative ring with a unique maximal ideal. A Dedekind domain with a unique maximal ideal, so a local integrally closed, Noetherian domain that is not a field, is called a **discrete valuation ring**. If $A$ is a Dedekind domain and $\mathfrak{p}$ a maximal...
ideal in \( A \), then the localization \( A_p \) of \( A \) is a discrete valuation ring, e.g., \( \mathbb{Z}_{(p)} \) is a discrete valuation ring where \( p \) is a prime in \( \mathbb{Z} \) as is \( F[t]_{(f)} \) with \( F \) a field and \( f \) an irreducible polynomial in \( F[t] \).

Remark 78.4. If \( A \) is a nonzero commutative ring, we define the (Krull) dimension of \( A \) by

\[
\dim A = \max \{ n \mid \text{there exists a chain of prime ideals } p_0 < \cdots < p_n \text{ in } A \},
\]

if this a finite number (and to be infinite if not). (The trivial ring is said to have (Krull) dimension \(-1\) (or \(-\infty\).) With this definition, a Dedekind domain is a Noetherian integrally closed domain of dimension one. The geometric analogue of a Dedekind domain turns out to be a smooth affine curve. The points of the curve correspond to the maximal ideals in a Dedekind domain.

Our extension theory allows us to conclude:

**Theorem 78.5.** Let \( A \) be a Dedekind domain with quotient field \( F \) and \( K/F \) a finite separable field extension. Then \( A_K \) is a Dedekind domain.

**Proof.** \( A_K \) is integrally closed by definition and Noetherian by Theorem 77.5. So it suffices to show if \( \mathfrak{P} \) is a nonzero prime ideal in \( A_K \), then it is a maximal ideal, equivalently, \( A_K/\mathfrak{P} \) is a field. By the lemma, \( \mathfrak{P} \cap A \) is a prime ideal. If \( x \) is a nonzero element in \( \mathfrak{P} \), it satisfies an equation

\[
(*) \quad x^n + a_{n-1}x^{n-1} + \cdots + a_0 = 0 \text{ for some } a_0, \ldots, a_{n-1} \text{ in } A.
\]

As \( A_K \) is a domain, we may assume that \( a_0 \) is nonzero. Since \( a_0 \) lies in \( \mathfrak{P} \cap A \), the prime ideal \( p = \mathfrak{P} \cap A \) is nonzero hence maximal, so \( A/p \) is a field. But \( A_K \) is a finitely generated \( A \)-module by Theorem 77.5, so the domain \( A_K/\mathfrak{P} \) is a finite dimensional \((A/p)\)-vector space. It follows that \( A_K/\mathfrak{P} \) is a field by Corollary 47.18. \( \square \)

Using equation \((*)\) above, we easily see that the following is true:

**Corollary 78.6.** Let \( A \) be a domain with quotient field \( F \) and \( K/F \) a finite field extension. If \( \mathfrak{A} \) is a nonzero ideal in \( A_K \), then \( \mathfrak{A} \cap A_F \) is nonzero.

**Corollary 78.7.** Let \( A \) be a PID with quotient field \( F \) and \( K/F \) a finite separable extension. Then \( A_K \) is a Dedekind domain. Moreover, \( A_K \) has an integral basis.

**Remarks 78.8.** 1. The corollary shows that every ring of algebraic integers is a Dedekind domain and is a finitely generated free abelian group.
2. The corollary shows that $F[t]_K$ for any field $F$ and finite separable extension $K$ of $F(t)$ is also a Dedekind domain and is a finitely generated free $F[t]$-module. It turns out that the theory about rings of algebraic integers and $F[t]_K$ is very similar when $F$ is a finite field. This case can be viewed as the geometric analogue of the arithmetic case of the ring of algebraic integers. Although one must worry about separability, it is easier in general. A field that is either a finite extension of the rationals or a finite separable extension of $F(t)$ with $F$ a finite field is called a global field. Since the expectation that study of the rings $\mathbb{Z}_K$ and $F[t]_K$ in the appropriate global field $K$ are similar, this often leads to the search and study for the appropriate analogue. For example, the solution of the analogue of the Riemann Hypothesis (the Weil Conjecture) for the geometric case by Deligne is one of the major theorems in mathematics proven in the twentieth century.

3. In general, if $A$ is a Dedekind domain with quotient field $F$ and $K/F$ a finite separable field extension, then $A_K$ is not a free $A$-module.

4. Let $A$ be a domain with quotient field $F$ and $K/F$ a finite field extension. In general, $A_K$ will not be a finitely generated $A$-module. However, the Krull-Akizuki Theorem 93.17 below says if $A$ is a Noetherian domain of dimension at most one, then $B$ has the same properties for any domain $B$ satisfying $A \subset B \subset K$. In particular, if $A$ is a Dedekind domain, then $A_K$ is also a Dedekind domain. In particular, if $K/F(t)$ is a finite extension, then $F[t]_K$ is a Dedekind domain. Therefore, any finite extension of $F(t)$ is also called a global field.

5. Let $A$ be a Dedekind domain with quotient field $F$ and $K/F$ a finite separable field extension. If $\mathfrak{p}$ is a maximal ideal in $A$, then $(A_\mathfrak{p})_K = (A_K)_\mathfrak{p}$, the localization at $A \setminus \mathfrak{p}$ by Exercise 76.15(5). In particular, $(A_\mathfrak{p})_K$ is a Dedekind domain. Although $A_\mathfrak{p}$ is a discrete valuation ring, $(A_\mathfrak{p})_K$ may not be. However, by Exercise 78.14 (6), it does have only finitely many maximal ideals and both it and $A_\mathfrak{p}$ are PIDs.

We next prove the defining property of Dedekind domains that also achieves our desired generalization of a UFD.

**Theorem 78.9.** Let $A$ be a Dedekind domain. Then every nonzero nonunit ideal in $A$ is a product of nonzero prime ideals, unique up to order.

**Proof.** Let $\mathfrak{a} < A$ be a nonzero ideal and $F$ the quotient field of $A$. We prove the result in a number of steps.
The first step uses the fact that $A$ is Noetherian.

**Step 1.** There exist nonzero prime ideals (hence maximal ideals) $p_1, \ldots, p_r$ in $A$, not necessarily distinct, satisfying $p_1 \cdots p_r \subseteq A$ (Cf. Exercise 29.22 (13)):

If this is not true, by the Maximal Principle, we may assume that $A$ is a maximal counterexample, i.e., $A$ does not satisfy Step 1, but any nonunit ideal properly containing $A$ does. Of course, $A$ cannot be a prime ideal, so there exist ideals $B_i$, $i = 1, 2$ satisfying $A < B_i < A$ with $B_1 B_2 \subseteq A$ by Lemma 25.20. As each $B_i$ contains a product of prime ideals by the maximality condition, so does $A$, a contradiction.

The second step uses the fact that nonzero prime ideals in the Noetherian domain $A$ are maximal.

**Step 2.** Let $p$ be a nonzero prime ideal in $A$. Set

$$p^{-1} := \{x \in F \mid xp \subseteq A\},$$

an $A$-submodule of $F$. Then $A < p^{-1}$:

Clearly, $A \subset p^{-1} \subset F$ and $p^{-1}$ is an $A$-module. Let $a$ be a nonzero element of $p$. If $p = (a)$, then $a^{-1}$ lies in $p^{-1} \setminus A$, so we may assume not. By Step 1, there exist nonzero prime ideals $p_1, \ldots, p_r$ satisfying $p_1 \cdots p_r \subseteq (a) \subset p < A$. Moreover, we may assume that we have chosen these prime ideals so that $r > 1$ is minimal. Since $p$ is a prime ideal, $p_i \subset p$ for some $i$, and as $p_i$ is a maximal ideal, we, in fact, must have $p_i = p$. Changing notation if necessary, we may assume that $p = p_1$.

By the minimality of $r$, there exists an element $b$ in $p_2 \cdots p_r \setminus (a)$. We then have

$$bp = bp_1 \subset p_1 \cdots p_r \subset (a).$$

In particular, $a^{-1}bp$ lies in $A$, i.e., $a^{-1}b$ lies in $p^{-1}$. By the choice of $b$, we have $a^{-1}b \notin A$.

The third step uses the fact that the dimension one Noetherian domain $A$ is integrally closed.

**Step 3.** Let $p$ be a nonzero prime ideal in $A$. Then for every nonzero ideal $\mathfrak{B}$ in $A$, we have $\mathfrak{B} < \mathfrak{B} p^{-1}$. In particular, $pp^{-1} = A$:

Suppose that $\mathfrak{B} = \mathfrak{B} p^{-1}$. By Step 2, there exists an element $a$ in $p^{-1} \setminus A$ satisfying $a \mathfrak{B} \subset \mathfrak{B}$. Hence $\lambda_a : \mathfrak{B} \to \mathfrak{B}$ is an $A$-homomorphism. Therefore, $\mathfrak{B}$ is an $A[a]$-module. As $A[a]$ is a nonzero ring, $\text{ann}_{A[a]} \mathfrak{B} = 0$ (as $\mathfrak{B} \subset F$ is $A[a]$-torsion free). Moreover, $\mathfrak{B}$ is finitely generated, since $A$ is a Noetherian ring. Therefore, $a$ is integral over $A$ by Claim 76.4. As $A$ is integrally closed, we conclude that $a$ lies in $A$, a contradiction. Finally, if $\mathfrak{B} = p$, then $p < pp^{-1}$ forces $pp^{-1}$ to be $A$, as $p$ is a maximal ideal.
Step 4. Finish:

We first show that \( \mathfrak{A} \) is a product of nonzero prime ideals. If not, as in the proof of Step 1, we may assume that \( \mathfrak{A} \) is a maximal counterexample. Since \( \mathfrak{A} \) cannot be a prime ideal, there exists a prime ideal \( \mathfrak{p} \) in \( A \) with \( \mathfrak{A} \subsetneq \mathfrak{p} \). In particular, \( \mathfrak{A}\mathfrak{p}^{-1} \subset \mathfrak{p}\mathfrak{p}^{-1} = A \). By Step 3 and the maximality condition, there exist nonzero prime ideals \( \mathfrak{p}_1, \ldots, \mathfrak{p}_r \) such that \( \mathfrak{A}\mathfrak{p}^{-1} = \mathfrak{p}_1 \cdots \mathfrak{p}_r \). Applying Step 4 again shows that \( \mathfrak{A} = \mathfrak{p}_1 \cdots \mathfrak{p}_r \), a contradiction.

Next we show uniqueness. If we have

\[ p_1 \cdots p_r = \mathfrak{A} = \mathfrak{p}_1 \cdots \mathfrak{p}_s \]

with all the \( p_i \) and \( \mathfrak{p}_j \) maximal ideals, then as in Step 1, we see that \( p_i = \mathfrak{p}_j \) for some \( i \) and \( j \). Changing notation, we may assume that \( 1 = i = j \). Multiplying the above equation by \( p^{-1} \) leads to our conclusion as \( p_1 p_1^{-1} = A \) by Step 3.

The converse of the theorem is true, i.e., a domain in which every nonzero nonunit ideal is a product of prime ideals, unique up to order, is a Dedekind domain. This is left as an exercise.

Let \( A \) be a Dedekind domain and \( \mathfrak{A} \) a nonzero nonunit ideal in \( A \). By the theorem, there exist unique nonzero prime ideals \( \mathfrak{p}_1, \ldots, \mathfrak{p}_r \) in \( A \) and unique positive integers \( e_1, \ldots, e_r \) satisfying

\[ \mathfrak{A} = \mathfrak{p}_1^{e_1} \cdots \mathfrak{p}_r^{e_r} \]

(unique up to order). We call (*) a factorization of \( \mathfrak{A} \). For each nonzero prime ideal \( \mathfrak{p} \) and we define

\[ v_{\mathfrak{p}}(\mathfrak{A}) = \begin{cases} e_i & \text{if } \mathfrak{p} = \mathfrak{p}_i, \text{ some } i = 1, \ldots, r \\ 0 & \text{otherwise.} \end{cases} \]

Then (*) can be written

\[ \mathfrak{A} = \prod_{\mathfrak{p}} \mathfrak{p}^{v_{\mathfrak{p}}(\mathfrak{A})}. \]

Let \( \mathfrak{A} \) and \( \mathfrak{B} \) be two nonzero ideals in a Dedekind domain \( A \). Then, by unique factorization of ideals, \( \mathfrak{A} \subset \mathfrak{B} \) if and only if \( v_{\mathfrak{p}}(\mathfrak{A}) \geq v_{\mathfrak{p}}(\mathfrak{B}) \) for all maximal ideals \( \mathfrak{p} \) in \( A \). Mimicking the case for elements, we say \( \mathfrak{B} \) divides \( \mathfrak{A} \) and write \( \mathfrak{B} | \mathfrak{A} \). In particular, \( \mathfrak{p}^n | \mathfrak{A} \) if and only if \( n \leq v_{\mathfrak{p}}(\mathfrak{A}) \). We define the greatest common divisor of \( \mathfrak{A} \) and \( \mathfrak{B} \) in \( A \) to be the largest ideal \( \mathfrak{D} \) of \( A \) containing both \( \mathfrak{A} \) and \( \mathfrak{B} \). As \( \mathfrak{A} \) and \( \mathfrak{B} \) lie in \( \mathfrak{A} + \mathfrak{B} \), it follows that \( \mathfrak{A} + \mathfrak{B} = \mathfrak{D} = \prod_{\mathfrak{p}} \mathfrak{p}^{\min\left(v_{\mathfrak{p}}(\mathfrak{A}), v_{\mathfrak{p}}(\mathfrak{B})\right)} \). We say \( \mathfrak{A} \) and \( \mathfrak{B} \) are relatively prime if \( \mathfrak{D} = A \), i.e., \( A = \mathfrak{A} + \mathfrak{B} \).
We have shown that if $A$ is a Dedekind domain, then any nonzero prime ideal $p$ has an ‘inverse’, viz., $p^{-1}$ as $pp^{-1} = A$. We next generalize this to a notion of inverses of ideals in Dedekind domains.

**Definition 78.10.** Let $A$ be a domain with quotient field $F$. A nonzero $A$-submodule $\mathfrak{A}$ of $F$ is called a fractional ideal of $A$ if there exists a nonzero element $x$ in $A$ satisfying $x\mathfrak{A} \subset A$. If $\mathfrak{A}$ is a fractional ideal, we define

$$\mathfrak{A}^{-1} := \{ x \in F \mid x\mathfrak{A} \subset A \}.$$  

As this is clearly a submodule of $F$ and $a\mathfrak{A}^{-1} \subset A$ for any $a$ in $\mathfrak{A}$, the module $\mathfrak{A}^{-1}$ is also a fractional ideal. Set

$$I_A := \{ \mathfrak{A} \mid \mathfrak{A} \text{ a fractional ideal of } A \}.$$  

We say that $\mathfrak{A}$ in $I_A$ is invertible if $\mathfrak{A}\mathfrak{A}^{-1} = A$.

As every maximal ideal in a Dedekind domain $A$ is invertible and there exist unique factorization of nonzero nonunit ideals in $A$, we expect that every fractional ideal in $A$ is invertible. This is in fact true, as we shall now show.

**Corollary 78.11.** Let $A$ be a Dedekind domain. Then every fractional ideal of a Dedekind domain $A$ is invertible. Moreover, $I_A$ is a free abelian group on basis the set of maximal ideals in $A$.

**Proof.** Let $\mathfrak{A}$ be a fractional ideal. Then there exists a nonzero element $a$ in $A$ satisfying $a\mathfrak{A} \subset A$. By the theorem, we have factorizations say $a\mathfrak{A} = \mathfrak{P}_1^{e_1} \cdots \mathfrak{P}_r^{e_r}$ and $(a) = p_1^{f_1} \cdots p_s^{f_s}$. As nonzero prime ideals are invertible, we have $\mathfrak{A} = \mathfrak{P}_1^{e_1} \cdots \mathfrak{P}_r^{e_r} p_1^{-f_1} \cdots p_s^{-f_s}$. Cancelling appropriate $\mathfrak{P}_i$’s and $p_j$’s gives a factorization that is clearly unique. It follows that $I_A$ is a free abelian group on the set of maximal ideals with unity $A$. □

Of course, if $A$ is a Dedekind domain and $\mathfrak{A}$ a fractional ideal, we have $\mathfrak{A} = \prod_p p^{v(p^n)}$, where $v_p(p^n) = n$ for any integer $n$.

**Definition 78.12.** Let $A$ be a Dedekind domain. Set

$$P_A := \{ xA \mid x \in F^\times \}$$

a subgroup of $I_A$ called the group of principal fractional ideals of $A$. The quotient group

$$Cl_A = I_A/P_A$$

is called the class group of $A$.

It measures the obstruction of $A$ from being a PID. It also measures when $A$ is a UFD as we shall now show.
Corollary 78.13. Let $A$ be a Dedekind domain. Then $A$ is a UFD if and only if $A$ is a PID if and only if $Cl_A$ is trivial.

Proof. We know that a domain is a UFD if and only if every nonzero prime ideal contains a prime element by Kaplansky’s Theorem 30.1. As $A$ is a Dedekind domain, this is equivalent to every maximal ideal being principal which, by Theorem 78.9, is equivalent to every ideal being principal.

In general, a Dedekind domain may not be a PID, e.g., $\mathbb{Z}[\sqrt{-5}]$. Also the class group $Cl_A$ may be infinite. In fact L. Claburn has shown that any abelian group can be realized as the class group of some Dedekind domain. However, an important theorem in number theory is that $Cl_{\mathbb{Z}_K}$ is a finite group for any ring of algebraic integers $\mathbb{Z}_K$. This is usually proven by using what is called Minkowski Theory.

Exercises 78.14.

1. Define the least common multiple of two nonzero nonunit ideals in a Dedekind domain and show that it is equal to the intersection of these ideals.

2. Show that $\mathbb{Z}[\sqrt{-5}]$ is a Dedekind domain but not a PID.

3. Let $A$ be a discrete valuation ring, i.e., a Dedekind domain with a unique maximal ideal. Show that $A$ is a PID.

4. Let $A$ be a Dedekind domain and $S$ a multiple set not containing zero. Show that the localization $S^{-1}A$ is either a field or a Dedekind domain. In particular, if $\mathfrak{p}$ is a maximal ideal in $A$ then the localization of $A$ at $\mathfrak{p}$, i.e., $A_{\mathfrak{p}} = S^{-1}A$ where $S = A \setminus \mathfrak{p}$, is a discrete valuation ring.

5. Show that every Dedekind domain with finitely many prime ideals is a PID.

6. Let $A$ be a Dedekind domain with quotient field $F$ and having finitely many prime ideals. Show if $K/F$ is a finite separable field extension, then $A_K$ is a Dedekind domain with finitely many prime ideals.

7. Let $R$ be an arbitrary domain: Show the following:
   (a) Any invertible fractional ideal in $R$ is finitely generated.
   (b) A product of fractional ideals in $R$ is invertible if and only if each of the fractional ideals is invertible.
   (c) If an ideal in $R$ is a product of invertible prime ideals, then the prime ideals are unique up to order.
8. Let $R$ be a domain in which every nonzero ideal is a product of prime ideals. Show that every nonzero prime ideal is invertible and maximal.

9. Show that a domain, not a field, is a Dedekind domain, if and only if every nonzero ideal is a product of prime ideals.

10. Let $A$ be a Dedekind domain and $\mathfrak{p}$ a prime ideal in $A$. Show the following:
   (i) $\mathfrak{p}^m = \mathfrak{p}^n$ if and only if $m = n$.
   (ii) If $\pi \in \mathfrak{p} \setminus \mathfrak{p}^2$, then $\mathfrak{p}^m = (\pi)^m + \mathfrak{p}^n$ for any positive integers $n \geq m$.
   (iii) The only proper ideals in $A/\mathfrak{p}^n$ are $\mathfrak{p}/\mathfrak{p}^n, \ldots, \mathfrak{p}^{n-1}/\mathfrak{p}^n$.

11. Let $A$ be a Dedekind domain and $\mathfrak{a}$ a nonzero ideal in $A$. Using the previous exercise show the following:
   (i) Every ideal in $A/\mathfrak{a}$ is principal.
   (ii) Every ideal in $A$ can be generated by two elements.

12. Let $A$ be a Dedekind domain. Prove the following:
   (i) Let $\mathfrak{A}$ be a fractional ideal of $A$ and $\mathfrak{B} \subset A$ an ideal. Then there exists an element $a \in \mathfrak{A}$ such that $\mathfrak{A}^{-1}a + \mathfrak{B} = A$.
   (ii) Let $\mathfrak{A}_1$ and $\mathfrak{A}_2$ be two fractional ideals of $A$. Then the $A$-module $\mathfrak{A}_1 \bigoplus \mathfrak{A}_2$ is isomorphic to $A \bigoplus A_1 \mathfrak{A}_2$.
   (iii) Let $\mathfrak{A}_1$ be a fractional ideal of $A$. Then there exists an $A$-isomorphism $\mathfrak{A}_1 \bigoplus A^{-1} \cong A^2$. In particular, every fractional ideal over $A$ is $A$-projective (cf. Exercises 38.12(10), (11), (13)).
   (iv) Every finitely generated torsion-free $A$-module is isomorphic to a finite direct sum of ideals in $A$. (Cf. Exercise 43.24(2)).
   (v) Let $M$ be a finitely generated torsion-free $A$-module. Then $M \cong A^n \bigoplus \mathfrak{A}$ for some $n$ and some ideal $\mathfrak{A}$ in $A$. [There is a uniqueness statement. Can you guess what it is?]

13. Let $A$ be a local noetherian domain of dimension one with $\mathfrak{m}$ its maximal ideal. Let $k := A/\mathfrak{m}$, called the residue class field of $A$. Show the following are equivalent:
   (i) $A$ is a discrete valuation ring.
   (ii) $A$ is integrally closed.
   (iii) $\mathfrak{m}$ is a principal ideal.
   (iv) $\dim_k(\mathfrak{m}/\mathfrak{m}^2) = 1$.
   (v) Every non-zero ideal is a power of $\mathfrak{m}$.
   (vi) There exists $x \in A$ such that every non-zero ideal is of the form $(x^n)$, $n \geq 0$. 

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(vii) $A$ is a valuation ring, i.e., if a nonzero $x$ is an element in the quotient field of $A$, then either $x$ is an element of $A$ or its multiplicative inverse is.

(viii) $A$ is a Dedekind domain.

14. Let $A$ be a noetherian domain. Then the following are equivalent:
   
   (i) $A$ is a Dedekind domain.
   
   (ii) $A_m$ is a discrete valuation ring for all maximal ideals $m$ in $A$.
   
   (iii) Every non-zero fractional ideal in $A$ is invertible.
   
   (iv) Every integral ideal in $A$ factors uniquely into prime ideals.

15. Let $A$ be a Dedekind domain and $\mathfrak{A}$ an nonzero ideal in $A$ with factorization $\mathfrak{A} = p_1^{e_1} \cdots p_r^{e_r}$ into prime ideals in $A$. Prove all of the following:
   
   (i) There exists a ring isomorphism $A/\mathfrak{A} \cong A/p_1^{e_1} \times \cdots A/p_r^{e_r}$.
   
   (ii) If $M$ is an $A$-module, then there exist $A$-modules $M_i$ for $i = 1, \ldots, r$ and an $A$-module isomorphism $M \cong \prod_{i=1}^{r} M_i$ with $A/p_i^{k_j} M_i = 0$ for every $j \neq i$ and each $i = 1, \ldots, r$. In particular, each $M_i$ is an $A/p_i^{e_i}$-module.

16. Let $A$ be a Dedekind domain, $p$ a nonzero prime ideal in $A$, and $M$ a finitely generated torsion $A$-module with annihilator, $\text{ann}_A M := \{x \in A | xA = 0\} = \bigcap_{x \in A} \text{ann}_A x$. Show all of the following:
   
   (i) Let $A_p$ be the localization of $A$ at $p$, a discrete valuation ring hence a PID. Then $A/p^e \cong A_p/(p^e A_p)$ as rings.
   
   (ii) If $M$ satisfies $p^e = \text{ann}_A(M)$, then $M$ is isomorphic to a coproduct of $A$-modules $A/p^k$ for various $1 \leq k \leq e$.
   
   (iii) If $\text{ann}_A(M) = p_1^{e_1} \cdots p_r^{e_r}$, then $M$ is isomorphic to a coproduct of $A$-modules $A/p_i^{k_{ij}}$ for various $1 \leq k_{ij} \leq e_i$ for $i = 1, \ldots, r$. [There is a uniqueness statement. What is it?]

17. Let $A$ be a Dedekind domain and $M$ a finitely generated $A$-module. Then $M = M_{tf} \oplus M_t$ with the submodules $M_{tf}$ torsion-free and $M_t$ torsion, respectively. In particular, by (12) and (16) of these Exercises (78.14), $M \cong R^n \prod \mathfrak{A} \prod \bigoplus_{i,j} A/p_i^{k_{ij}}$, for some ideal $\mathfrak{A}$ in $A$ and for various $A/p_i^{k_{ij}}$, $1 \leq k_{ij} \leq e_i$, where $\text{ann}_A(M) = p_1^{e_1} \cdots p_r^{e_r}$. (There is also a uniqueness statement. What is it?)

79. Extension of Dedekind Domains

Let $A$ be a Dedekind domain with quotient field $F$ and $K/F$ a finite separable field extension. We have shown that $A_K$ is also a Dedekind domain. We wish to investigate how nonzero prime ideals in $A$ factor into a product of primes ideals in $A_K$. We begin by showing any proper ideal in $A$ remains a proper ideal in $A_K$. This follows from the following
Claim 79.1. Let \( p \) be a nonzero prime ideal in the Dedekind domain \( A \) above. Then \( pA_L \not< A_L \) for any field extension \( L/F \).

We know that \( p^2 < p \). Let \( \pi \in p \setminus p^2 \). It follows that \( \pi A = pA \) for some ideal \( A \) of \( A \) with \( p \not\mid A \), hence \( p + A = A \). Write \( 1 = p + a \) with \( p \in p \) and \( a \in A \). It follows that \( a \not\in p \) and \( ap \subseteq pA = \pi A \). Suppose that \( pA_L = A_L \). Then \( aA_L = apA_L \subseteq \pi A_L \), hence \( a = \pi x \), for some \( x \in A_L \cap F = A \). This implies that \( a \) lies in \( p \), a contradiction. This proves the claim.

In particular, if \( p \) is a nonzero prime ideal in \( A \), then we have a factorization
\[
(*) \quad pA_K = \mathfrak{P}_1^{e_1} \cdots \mathfrak{P}_r^{e_r} \quad \text{in} \quad A_K
\]
called the splitting behavior of \( p \) in \( A_K \). We say a prime ideal \( \mathfrak{P} \) in \( A_K \) lies over \( p \), if \( \mathfrak{P} \cap A = p \). If \( \mathfrak{P} \) is such a prime ideal then \( pA_K \subseteq (\mathfrak{P} \cap A)A_K \subseteq \mathfrak{P} < A_K \), so \( \mathfrak{P} = \mathfrak{P}_i \) for some \( i \), as \( \mathfrak{P} \) is a maximal ideal in \( A_K \). Conversely if \( \mathfrak{P} = \mathfrak{P}_i \), then \( pA_K \cap A \subset \mathfrak{P} \cap A < A \), so \( \mathfrak{P} \cap A = p \), as \( p \) is a maximal ideal. Thus
\[
\{ \mathfrak{P}_1, \ldots, \mathfrak{P}_r \} = \{ \mathfrak{P} \mid \text{with} \ \mathfrak{P} \ \text{a prime ideal in} \ A_K \ \text{satisfying} \ \mathfrak{P} \ | pA_K \},
\]
the set of prime ideals in \( A_K \) lying over \( p \) in \( A_K \). We shall write \( \mathfrak{P} | p \) for \( \mathfrak{P} | pA_K \). Let \( \mathfrak{P} \) be a nonzero prime ideal in \( A_K \). Then the composition of the inclusion \( A \subset A_K \) and the canonical epimorphism \( \gamma : A_K \to A_K/\mathfrak{P} \) shows that the field \( A_K/\mathfrak{P} \) is a field extension of \( A/\mathfrak{P} \cap A \). We define ramification index \( e(\mathfrak{P}/p) \) of \( \mathfrak{P} \) over \( p \) and if \( \mathfrak{P} \) lies over \( p \), the inertia index \( f(\mathfrak{P}/p) \) of \( \mathfrak{P} \) over \( p \) by
\[
e(\mathfrak{P}/p) = v_\mathfrak{P}(pA_K)
\]
\[
f(\mathfrak{P}/p) = [A_K/\mathfrak{P} : A/p].
\]
respectively. If \( e(\mathfrak{P}/p) > 1 \) or if \( (A_K/\mathfrak{P})/(A/p) \) is not separable, we say that \( p \) ramifies in \( A_K \). That the inertia index is a finite number follows from the following important result:

Theorem 79.2. Let \( A \) be a Dedekind domain with quotient field \( F \), \( K/F \) a finite separable field extension. If \( p \) is a nonzero prime ideal in \( A \), then \( f(\mathfrak{P}/p) \) is finite for all \( \mathfrak{P} | p \) and
\[
[K : F] = \sum_{\mathfrak{P} | p} e(\mathfrak{P}/p)f(\mathfrak{P}/p).
\]
Proof. Let $pA_K = \mathfrak{p}_1^{e_1} \cdots \mathfrak{p}_r^{e_r}$ be a factorization of $p$ in $A_K$. So $e(\mathfrak{p}_i/p) = e_i$ and $f(\mathfrak{p}_i/p) = f_i$ for each $i$. Let $n = [K : F]$ and $\mathfrak{A} : A_K \to A_K/pA_K$ be the canonical epimorphism, so $\mathfrak{A} = A/p$. By the Chinese Remainder Theorem (*),

$$A_K/pA_K = A_K/\mathfrak{p}_1^{e_1} \times \cdots \times A_K/\mathfrak{p}_r^{e_r}.$$ 

The ring $\mathfrak{A}_K = A_K/pA_K$ is a finite dimensional $\mathfrak{A}$-vector space as $A_K$ is a finitely generated $A$-module, hence so are the $A_K/\mathfrak{p}_i^{e_i}$'s. It follows from this that the $f(\mathfrak{p}_i/p)$ are all finite.

**Step 1.** $n = [K : F] = \dim_{\mathfrak{A}} \mathfrak{A}_K$: 

Choose $B = \{x_1, \ldots, x_m\}$ in $A_K$, so that $B = \{\mathfrak{x}_1, \ldots, \mathfrak{x}_m\}$ is an $\mathfrak{A}$-basis for $\mathfrak{A}_K$.

We shall show that $B$ is an $F$-basis for $K$ establishing the claim. We first show $B$ is linearly independent. Suppose not, then clearing denominators if necessary, we see that we have an equation

$$a_1x_1 + \cdots + a_mx_m = 0 \quad \text{in} \quad A_K$$

for some $a_1, \ldots, a_m$ in $A$, not all zero. Let $\mathfrak{a}$ be the nonzero ideal $(a_1, \ldots, a_m)$. As $\mathfrak{a}\mathfrak{a}^{-1} = A$, we can choose a nonzero element $a$ in $\mathfrak{a}^{-1}$ such that not every $aa_i$ lies in $p$. As $aa_i \in A$ for every $i$, we have

$$\underbrace{aa_1x_1 + \cdots + aa_mx_m}_\text{in} \quad \mathfrak{A}_K$$

contradicting $B$ is $\mathfrak{A}$-linearly independent. Thus $B$ is linearly independent.

We next show $B$ spans. Let $B = Ax_1 + \cdots + Ax_m$, a submodule of $A_K$. As $B$ is a basis for $\mathfrak{A}_K$, if $b \in B$, there exists an $x \in A_K$ satisfying $x - b$ lies in $pA_K$, i.e., $A_K = B + pA_K$. Let $M = A_K/B$, a finitely generated $A$-module. Then $M/pM = 0$, since each $x_i$ maps to zero in it so a zero dimensional $\mathfrak{A}$-vector space. Consequently, $M = pM$. Set $M = Ay_1 + \cdots + Ay_s$. Then we have equations

$$y_i = \sum_{j=1}^s p_{ij}y_j \quad \text{for some} \quad p_{ij} \in p, \quad 1 \leq i, j \leq s,$$

hence

$$\sum_{j=1}^s (\delta_{ij} - p_{ij})y_j = 0 \quad \text{for} \quad i = 1, \ldots, s.$$ 

Let $\Delta$ be the determinant of the $s \times s$ matrix $(\delta_{ij} - p_{ij})$. Then $\Delta y_i = 0$ for every $i$ by Cramer’s Rule, so $\Delta M = 0$. It follows that $\Delta A_K \subset B$, 

hence also that $\Delta K \subset Fx_1 + \cdots + Fx_m$. But $\Delta \equiv 1 \mod p$, so $\Delta$ is nonzero. Therefore, we conclude that

$$K = \Delta K = Fx_1 + \cdots + Fx_m.$$  

Therefore, $\mathfrak{B}$ is an $F$-basis for $K$ and $m = n$. This completes Step 1.

**Step 2.** Finish:

In view of (*), it suffices to show that $\dim_{\overline{A}} A_K/\mathfrak{P}_i^{e_i} = e_if_i$ for each $i$, $1 \leq i \leq r$.

We have a descending chain of $\overline{A}$-vector spaces

$$A_K/\mathfrak{P}_i^{e_i} \supset \mathfrak{P}_i/\mathfrak{P}_i^{e_i} \supset \cdots \supset \mathfrak{P}_i^{e_i-1}/\mathfrak{P}_i^{e_i} \supset 0,$$

so

$$\dim_{\overline{A}} A_K/\mathfrak{P}_i^{e_i} = \sum_{m=0}^{e_i-1} \dim_{\overline{A}} \mathfrak{P}_i^m/\mathfrak{P}_i^{m+1}.$$  

(Why?) As $\mathfrak{P}_i^{m+1} \subset \mathfrak{P}_i^m$ by unique factorization of ideals, we can choose an element $a$ in $\mathfrak{P}_i^m \setminus \mathfrak{P}_i^{m+1}$ for each $m$. Then $\mathfrak{P}_i^m = aA_K + \mathfrak{P}_i^{m+1}$, the greatest common divisor of $\mathfrak{P}_i^{m+1}$ and $aA_K$ and $\mathfrak{P}_i^{m+1} = aA_K \cap \mathfrak{P}_i^{m+1}$ (the least common multiple of $\mathfrak{P}_i^{m+1}$ and $aA_K$, cf. Exercise 78.14(1)). Let $\lambda_a : A_K \to \mathfrak{P}_i^m/\mathfrak{P}_i^{m+1}$ be the homomorphism defined by $x \mapsto ax + \mathfrak{P}_i^{m+1}$. The kernel of $\lambda_a$ is $\mathfrak{P}_i$ and $\lambda_a$ is onto, as

$$\text{im } \lambda_a = (aA_K + \mathfrak{P}_i^{m+1})/\mathfrak{P}_i^{m+1} = \mathfrak{P}_i^m/(aA_K \cap \mathfrak{P}_i^{m+1}) = \mathfrak{P}_i^m/\mathfrak{P}_i^{m+1}.$$  

Therefore $f_i = \dim_{\overline{A}} A_K/\mathfrak{P}_i = \dim_{\overline{A}} \mathfrak{P}_i^m/\mathfrak{P}_i^{m+1}$ for any $m > 0$ and $\dim_{\overline{A}} A_K/\mathfrak{P}_i^{e_i} = e_if_i$ as desired. \hfill \square

**Remark 79.3.** If $\mathfrak{P} | \mathfrak{p}$ in the theorem and $(A_K/\mathfrak{P})_{\text{sep}}$ is the separable closure of $A/\mathfrak{p}$ in $A_K/\mathfrak{P}$, then $f(\mathfrak{P}/\mathfrak{p}) = [(A_K/\mathfrak{P})_{\text{sep}} : A/\mathfrak{p}]p^s$ for some $s$ where $p = \text{char}(A/\mathfrak{p})$. In particular, if $(A_K/\mathfrak{P})/(A/\mathfrak{p})$ is normal, then $f(\mathfrak{P}/\mathfrak{p}) = |G(A_K/\mathfrak{P}]/(A/\mathfrak{p})|p^s$, so if $A_K/\mathfrak{P}$ is also separable, then $f(\mathfrak{P}/\mathfrak{p}) = |G((A_K/\mathfrak{P}]/(A/\mathfrak{p}))|$.

**Remark 79.4.** If $A = \mathbb{Z}$ in the theorem, Step 1 can be simplified. Indeed let $\{w_1, \ldots, w_n\}$ be an integral basis for $\mathbb{Z}_K$ and $\mathfrak{A}$ any nonzero ideal in $\mathbb{Z}_K$. We know that $\mathfrak{A} \cap \mathbb{Z} = (a)$ for some nonzero integer $a$. Set $S = \{\sum r_iw_i \mid 0 \leq r_i < a\}$. Then using the division algorithm, one shows that $S$ is a system of representatives for the cosets $\mathbb{Z}_K/\mathfrak{A}$, so has $a^n$ elements.

Let $A$ be a Dedekind domain with quotient field $F$ and $K/F$ a finite separable extension. Let $\mathfrak{p}$ be a nonzero prime ideal in $A$ and $\mathfrak{P}$ a prime ideal in $A_K$ lying over $\mathfrak{p}$. We say that $\mathfrak{P}$ is *unramified* over $A$ if $e(\mathfrak{P}/\mathfrak{p}) = 1$ and $A_k/\mathfrak{P}$ is a separable extension of $A/\mathfrak{p}$ and *ramified* otherwise (and *totally ramified* if in addition $f(\mathfrak{P}/\mathfrak{p}) = 1$). We say $\mathfrak{p}$ is
unramified in $A_K$ if every prime ideal in $A_K$ lying over $p$ is unramified and the extension $K/F$ (or $A_K/A$) is called unramified if every prime ideal in $A$ is unramified in $A_K$. We say a nonzero prime ideal $p$ in $A$ splits completely in $A_K$ if $e(P/p) = f(P/p) = 1$ for every prime ideal $P$ lying over $p$. If $K$ is a global field, with $A = \mathbb{Z}$ or $F[t]$, we do not have to worry about the separability condition of $A_K/P$ over $A/p$ as either $A_K/P$ is of characteristic zero or a finite field. A deep result in number theory states that there exist infinitely many nonzero prime ideals in $A$ that split completely in $A_K$ when $K$ is a global field.

The Kummer-Dedekind Theorem gives a method for computing the splitting behavior of prime ideals for all but finitely many prime ideals in a given separable extension. We first define the set of prime ideals whose splitting behavior this theorem omits. Let $A$ be a Dedekind domain with quotient field $F$ and $K/F$ a finite separable extension. By the Primitive Element Theorem 56.9, $K = F(\alpha)$ for some $\alpha \in K$. We may assume that $\alpha$ lies in $A_K$. So $A[\alpha]$ is a subring of $A_K$. Define the conductor of $A[\alpha]$ in $A_K$ to be the largest ideal $f$ of $A_K$ that lies in $A[\alpha]$, i.e.,

$$f = \{ x \in A_K \mid xA_K \subset A[\alpha] \}.$$ 

We know that $A_K$ contains the $F$-basis $\{1, \alpha_1, \ldots, \alpha^{n-1}\}$ for $K$, so it follows by Claim 77.6 that $f \neq 0$. Note that $f \subset A[\alpha]$ as $1$ lies in $A_K$, so the conductor $f$ is an ideal of both $A_K$ and $A[\alpha]$.

Kummer-Dedekind’s result is the following:

**Theorem 79.5.** (Kummer-Dedekind) Let $A$ be a Dedekind domain with quotient field $F$ and $p$ a nonzero prime ideal in $A$. Suppose that $K/F$ is a finite separable field extension with $K = F(\alpha)$, $\alpha \in A_K$, and $p$ relatively prime to the conductor $f$ of $A[\alpha]$ in $A_K$. Let $\tilde{\gamma} : A[t] \to (A/p)[t]$ be the canonical epimorphism and let $p_1, \ldots, p_r$ be distinct monic polynomials in $A[t]$ such that

$$m_F(\alpha) = p_1^{e_1} \cdots p_r^{e_r}$$

is a factorization of the image of the minimal polynomial of $\alpha$ in $(A/p)[t]$ into irreducibles. Set

$$P_i = pA_K + p_i(\alpha)A_K \text{ for each } i = 1, \ldots, r.$$ 

Then $P_1, \ldots, P_r$ are all of the distinct primes in $A_K$ lying over $p$, $f(P_i/p) = f_i$, and $e_i = e(P_i/p)$ for $i = 1, \ldots, r$. In particular,

$$pA_K = P_1^{e_1} \cdots P_r^{e_r}$$

is a factorization of $p$ in $A_K$. 
**Proof.** **Step 1.** \( A_K/pA_K \cong A[\alpha]/pA[\alpha] \):

Since \( pA_K \) and \( f \) are relatively prime, \( A_K = pA_K + f \). But \( f \subseteq A[\alpha] \), so we have \( A_K = pA_K + A[\alpha] \) and the natural map \( A[\alpha] \to A_K/pA_K \) is an epimorphism with kernel \( pA_K \cap A[\alpha] \). As \( f \cap A \) and \( p \) must also be relatively prime, we have \( pA[\alpha] + f = A[\alpha] \), so

\[
pA_K \cap A[\alpha] = (pA[\alpha] + f)(pA_K \cap A[\alpha]) \subset pA[\alpha].
\]

It follows that \( pA_K \cap A[\alpha] = pA[\alpha] \) establishing Step 1.

**Step 2.** \( \overline{A}[t]/(m_F(\alpha)) \cong A[\alpha]/pA[\alpha] \):

The kernel of the natural epimorphism \( A[t] \to \overline{A}[t]/(m_F(\alpha)) \) is generated by \( p \) and \((m_F(\alpha))\), Therefore, the evaluation map \( A[t] \to A[\alpha] \) induces the desired isomorphism as \( A[t]/pA[t] \cong (A/p)[t] \).

**Step 3.** Finish:

Let \( B = \overline{A}[t]/(m_F(\alpha)) \) and \( \sim : \overline{A}[t] \to B \) be the canonical epimorphism. By the Chinese Remainder Theorem,

\[
B \cong \overline{A}[t]/(\overline{p_i}^{e_i}) \times \cdots \times \overline{A}[t]/(\overline{p_r}^{e_r}).
\]

It follows that the nonzero prime ideals in \( B \) are the distinct principal ideals \((\overline{p_i})\) generated by the \( p_i \) mod \( m_F(\alpha) \), i.e., the \((p_i)\). We have \( \deg p_i = \deg \overline{p_i} = [B/(\overline{p_i}) : A] \) for \( i = 1, \ldots, r \). Further,

\[
0 = (m_F(\alpha)) = \bigcap_{i=1}^r (\overline{p_i}^{e_i}) \text{ in } B.
\]

By Steps 1 and 2, we know that \( B \cong A_K/pA_K \), with the isomorphism induced by evaluation at \( \alpha \), so by the Correspondence Principle, the prime ideals in \( B \) correspond to the prime ideals in \( A_K/pA_K \). Let \( \Omega_i \) in \( A_K/pA_K \) be the prime ideal corresponding to the prime ideal \((\overline{p_i})\) in \( B \) for \( i = 1, \ldots, r \). It follows that each \( \Omega_i \) is a principal ideal generated by \( \overline{p_i}(\alpha) \) with \( \deg \overline{p_i} = [A_K/\Omega_i : A/p] \) for each \( i \), and furthermore, 

\[
0 = \bigcap_{i=1}^r \Omega_i^{e_i}.
\]

Let \( \Psi_i \) be the ideal in \( A_K \) that is the preimage of \( \Omega_i \) under the natural epimorphism \( A_K \to A_K/pA_K \). Then \( \Psi_i \) is a nonzero prime ideal by Lemma 78.1 and satisfies \( \Psi_i = pA_K + p_i(\alpha)A_K \) for \( i = 1, \ldots, r \). It follows that \( \Psi_1, \ldots, \Psi_r \) are all the prime ideals in \( A_K \) lying over \( p \), and \( f_i = [A_K/\Psi_i : A/p] = \deg \overline{p_i} = \deg p_i \) for each \( i \). Moreover, as 

\[
e_i = |\{\Omega_i^m \mid m \geq 0\}|,
\]

so \( \Omega_i^{e_i} \) is the preimage of \( \Omega_i^{e_i} \). Since \( \prod_{i=1}^r \Psi_i^{e_i} \subset \bigcap_{i=1}^r \Psi_i^{e_i} \subset pA_K \), we have \( \nu_p(\Psi_i) \leq e_i \). We have \( [K : F] = \deg m_F(\alpha) = \sum_{i=1}^r e_i \deg p_i = \sum_{i=1}^r e_if_i \). If \( pA_K = \Psi_1^{e_1} \cdots \Psi_r^{e_r} \) is a factorization of \( pA_K \), then \( \sum_{i=1}^r e_if_i = \sum_{i=1}^r e_if_i \).
Let $K = \mathbb{Q}(\alpha)$ with $\alpha^3 = 2$. It can be shown that $\mathbb{Z}_K = \mathbb{Q}[\alpha]$. Then $t - 2 = (t - 3)(t^2 + 3t - 1)$ in $\mathbb{Z}/5\mathbb{Z}[t]$. As $t^2 + 3t - 1$ is irreducible in $\mathbb{Z}/5\mathbb{Z}[t]$, then by the Kummer-Dedekind Theorem 79.5, we see that $5\mathbb{Z}_K$ splits completely, say as $5\mathbb{Z}_K = p_1p_2$ with $f(p_1/5\mathbb{Z}_K) = 1$ and $f(p_2/5\mathbb{Z}_K) = 2$.

**Exercises 79.7.**

1. Complete the proof of Remark 79.4.
2. Let $A$ be a Dedekind domain with quotient field $F$ and $L/K/F$ finite separable extensions. If $\mathfrak{p}$ is a nonzero prime ideal in $A_L$, show

$$e(\mathfrak{P}/\mathfrak{P} \cap A) = e(\mathfrak{P}/\mathfrak{P} \cap A_K)e(\mathfrak{P} \cap A_K/\mathfrak{P} \cap A)$$

$$f(\mathfrak{P}/\mathfrak{P} \cap A) = f(\mathfrak{P}/\mathfrak{P} \cap A_K)f(\mathfrak{P} \cap A_K/\mathfrak{P} \cap A).$$

3. (Nakayama’s Lemma) Let $A$ be a local ring, i.e., a nonzero commutative ring with a unique maximal ideal, say $\mathfrak{m}$. Show if $M$ is a finitely generated $A$-module satisfying $\mathfrak{m}M = M$, then $M = 0$.

4. Let $A$ be a local ring with maximal ideal $\mathfrak{m}$ and $\mathfrak{m} : M \to M/\mathfrak{m}M$ the canonical epimorphism. Let $\mathcal{C} = \{x_1, \ldots, x_n\} \subset M$ and $\overline{\mathcal{C}} = \{\overline{x_1}, \ldots, \overline{x_n}\}$. Show that $\mathcal{C}$ generates $M$ if and only if $\overline{\mathcal{C}}$ spans $\overline{M}$ as an $A/\mathfrak{m}$-vector space and is a minimal generating set for $M$ (obvious definition) if and only if $\overline{\mathcal{C}}$ is an $A/\mathfrak{m}$-basis for $\overline{M}$.

5. Let $A$ be a Dedekind domain with quotient field $F$ and $K/F$ a finite separable field extension. If $\alpha$ is an element of $A_K$, then for each maximal ideal not dividing the conductor $f$ of $A[\alpha]$ (hence all but finitely many maximal ideals) $(A_{\mathfrak{p}})_K = A_{\mathfrak{p}}[\alpha]$, where $A_{\mathfrak{p}}$ is the localization of $A$ at $A \setminus \mathfrak{p}$.

6. Let $A$ be a Dedekind domain with quotient field $F$ and $K/F$ a finite separable extension. Then a prime ideal $\mathfrak{p}$ in $A$ ramifies in $A_K$ if and only if the localization $A_{\mathfrak{p}}$ of $A$ ramifies in $(A_{\mathfrak{p}})_K$.

**80. Hilbert Ramification Theory**

Theorem 79.2 becomes even nicer when $K/F$ is Galois. We need the following to prove this.

**Proposition 80.1.** Let $A$ be a Dedekind domain with quotient field $F$ and $\mathfrak{p}$ a nonzero prime ideal in $A$. If $K/F$ is a finite Galois extension, then the Galois group $G(K/F)$ acts transitively on the set of prime ideals in $A_K$ lying over $\mathfrak{p}$.
Theorem 79.2 allows us to use Galois Theory in the subject. This is called Hilbert Ramification Theory. We introduce the subject without going into it deeply (leaving some results as exercises).

For the rest of this section, let $A$ be a Dedekind domain with quotient field $F$ and $K/F$ a finite Galois extension. Let $p$ be a nonzero prime ideal in $A$ and $\mathfrak{P}$ a prime ideal in $A_K$ lying over $p$. As $G(K/F)$ acts on the set of prime ideals in $A_K$ lying over $p$, we can look at the isotropy subgroup

$$G_{\mathfrak{P}} = \{ \sigma \in G(K/F) \mid \sigma(\mathfrak{P}) = \mathfrak{P} \}$$

of $\mathfrak{P}$. It is called the decomposition group of $\mathfrak{P}$ and its fixed field

$$Z_{\mathfrak{P}} = \{ x \in K \mid \sigma(x) = x \text{ for all } \sigma \in G_{\mathfrak{P}} \}$$

is called the decomposition field of $\mathfrak{P}$. As usual, we have $\sigma G_{\mathfrak{P}} \sigma^{-1} = G_{\sigma(\mathfrak{P})}$. We also have the index $[G : G_{\mathfrak{P}}] = \text{number of primes in } A_K \text{ lying over } p$ as $G(K/F)$ acts transitively on the primes in $A_K$ lying over $p$. This means that $p$ splits completely in $A_K$ if and only if $G_{\mathfrak{P}} = 1$ and $G(K/F) = G_{\mathfrak{P}}$ if and only if $\mathfrak{P}$ is the only prime ideal in $A_K$ lying over $p$. 

**Proof.** Let $\mathfrak{P}_i \mid p$ be prime ideals in $A_K$ for $i = 1, 2$. Suppose that $\mathfrak{P}_2 \neq \sigma(\mathfrak{P}_1)$ for any $\sigma \in G(K/F)$. By the Chinese Remainder Theorem, there exists an element $x$ in $A_K$ satisfying $x \equiv 0 \mod \mathfrak{P}_2$ and $x \equiv 1 \mod \sigma(\mathfrak{P}_1)$ for all $\sigma \in G(K/F)$, i.e., $x \in \mathfrak{P}_2$ but $\sigma^{-1}(x) \notin \mathfrak{P}_1$ for any $\sigma$ in the group $G(K/F)$. Since

$$N_{K/F}(x) = \prod_{\sigma \in G(K/F)} \sigma(x) \text{ lies in } \mathfrak{P}_2 \cap A_F = p = \mathfrak{P}_1 \cap A \subset \mathfrak{P}_1,$$

we have $\sigma(x) \in \mathfrak{P}_1$ for some $\sigma \in G(K/F)$, a contradiction. □

**Proposition 80.2.** Let $A$ be a Dedekind domain with $F$ its quotient field and $p$ a nonzero prime ideal in $A$. If $K/F$ is a finite Galois extension, then $e(\mathfrak{P}/p) = e(\mathfrak{P}'/p)$ and $f(\mathfrak{P}/p) = f(\mathfrak{P}'/p)$ for all prime ideals $\mathfrak{P}$ and $\mathfrak{P}'$ lying over $p$ in $A_K$. In particular, if $\mathfrak{P} \mid p$ in $A_K$ is a prime ideal and $e = e(\mathfrak{P}/p)$, $f = f(\mathfrak{P}/p)$, then $[K : F] = efr$ where $r$ is the number of primes in $A_K$ lying over $p$.

**Proof.** By the last proposition, there exists an element $\sigma$ in $G(K/F)$ satisfying $\mathfrak{P}' = \sigma(\mathfrak{P})$, hence we have an isomorphism $A_K/\mathfrak{P} \to A_K/\sigma(\mathfrak{P})$ given by $x + \mathfrak{P} \mapsto \sigma(x) + \sigma(\mathfrak{P})$. It follows that $f(\mathfrak{P}/p) = f(\mathfrak{P}'/p)$. Further, if $\mathfrak{P}^m \mid pA_K$, then $\sigma(\mathfrak{P})^m \mid pA_K$, so $e(\mathfrak{P}/p) = e(\mathfrak{P}'/p)$. The result now follows using Theorem 79.2. □

**Proposition 80.2** allows us to use Galois Theory in the subject. This is called Hilbert Ramification Theory. We introduce the subject without going into it deeply (leaving some results as exercises).
Proposition 80.3. Let \( A \) be a Dedekind domain with quotient field \( F \) and \( K/F \) a finite Galois extension, \( \mathfrak{P} \) a prime ideal in \( A_K \) with \( \mathfrak{P} \) lying over \( p \). Then

1. \( \mathfrak{P} \) is the only prime ideal over \( \mathfrak{P} \cap Z_p \).
2. \( e(\mathfrak{P}/\mathfrak{P} \cap Z_p) = e(\mathfrak{P}/p) \) and \( f(\mathfrak{P}/\mathfrak{P} \cap Z_p) = f(\mathfrak{P}/p) \).
3. \( e(\mathfrak{P} \cap Z_p/p) = 1 \) and \( f(\mathfrak{P} \cap Z_p/p) = 1 \).

**Proof.** By Proposition 80.2, we have \([K : Z_p] = e' f' r'\) where \( e' \) is the ramification index of each prime ideal in \( A_K \) lying over \( \mathfrak{P} \cap Z_p \), \( f' \) its inertia index, and \( r' \) is the number of prime ideals in \( A_K \) lying over \( \mathfrak{P} \cap Z_p \). Similarly, we have \([Z_p : F] = e'' f'' r''\) where \( e'' \) is the ramification index of each prime ideal in \( Z_p \) over \( p \), \( f'' \) its inertia index, and \( r'' \) is the number of prime ideals in \( A_{Z_p} \) lying over \( p \), and \([K : F] = e f r\) where \( e \) is the ramification index of each prime ideal in \( A_K \) lying over \( p \), \( f \) is its inertia index, and \( r = [G(K/F) : G_{\mathfrak{P}}] \) is the number of prime ideals in \( A_K \) lying over \( p \). By the transitivity of ramification and inertia indices (cf. Exercise 79.7(2)), we have \( e = e' e'' \), \( f' f'' = f \), hence \( r = r' r'' \).

(1): By definition, \( G_{\mathfrak{P}} = G(K/Z_p) \), so the primes lying over \( \mathfrak{P}_{Z_p} \), \( A_{Z_p} \) by Proposition 80.1, each is \( \mathfrak{P} \) and \( r' = 1 \).

(2), (3): As \( r = [G(K/F) : G_{\mathfrak{P}}] = [Z_p : F] = e'' f'' r'' \), we have \([K : Z_p] = |G_{\mathfrak{P}}| = e f \). By (1), we have, \( e' e'' f' f'' = e f = [K : Z_p] = e' f' \). Thus \( e'' = 1 = f'' \) and \( e' = e \) and \( f' = f \).

We next relate the decomposition group \( G_{\mathfrak{P}} \) of \( \mathfrak{P} \) in \( A_K \) lying over \( p \) and the Galois group of the field extension \((A_K/\mathfrak{P})/(A/p)) \). Let \( \sigma \) be an element of the \( G_{\mathfrak{P}} \). Then \( \sigma \) induces an \((A/p)\)-isomorphism

\[ \sigma : A_K/\mathfrak{P} \rightarrow A_K/\mathfrak{P} \text{ given by } x + \mathfrak{P} \mapsto \sigma(x) + \mathfrak{P} = \sigma(x) + \mathfrak{P}. \]

Therefore, we get a group homomorphism

\[ (80.4) \quad - : G_{\mathfrak{P}} \rightarrow G((A_K/\mathfrak{P})/(A/p)) \text{ given by } \sigma \mapsto \sigma. \]

**Proposition 80.5.** Let \( A \) be a Dedekind domain with \( F \) its quotient field and \( p \) a nonzero prime ideal in \( A \). If \( K/F \) is a finite Galois extension and \( \mathfrak{P} \) a prime in \( A_K \) lying over \( p \), then \((A_K/\mathfrak{P})/(A/p)\) is a finite normal extension and the map \(- : G_{\mathfrak{P}} \rightarrow G((A_K/\mathfrak{P})/(A/p))\) in (80.4) is surjective.

**Proof.** Let \( - : A_K[t] \rightarrow (A_K/\mathfrak{P})[t] \) be the natural epimorphism. Since \( f(\mathfrak{P} \cap A_{Z_p}/p) = 1 \), we have \( A/p = A_{Z_p}/(\mathfrak{P} \cap A_{Z_p}) \), so we may assume that \( F = Z_p \), hence \( G(K/F) = G_{\mathfrak{P}} \). Let \( x \in A_K/\mathfrak{P} \) and choose \( \alpha \) in \( A_K \) satisfying \( \sigma x = x \). Then we have

\[ (*) \quad m_{A/p}(x) \mid m_F(\alpha) \text{ in } (A/p)[t]. \]
The kernel of (80.4) is denoted by $G_\mathfrak{P}$ in $A_{\mathfrak{P}}$. We may assume that $\alpha$ lies in $A_{\mathfrak{P}}$. Suppose that $\sigma$ lies in $G((A_{\mathfrak{P}})/(A/\mathfrak{p}))$. Then $\sigma$ is completely determined by $\sigma(\overline{\alpha})$. As $\sigma(\overline{\alpha})$ is a root of $m_{A/\mathfrak{p}}(\overline{\alpha})$, it is a root of $m_F(\alpha)$ by (*). Choose a root $\alpha'$ of $m_F(\alpha)$ satisfying $\alpha' = \sigma(\overline{\alpha})$. There exists an element $\sigma$ in $G(K/F)$ satisfying $\sigma(\alpha) = \alpha'$ as $G(K/F)$ acts transitively on the roots of $m_F(\alpha)$. Then $\sigma \mapsto \overline{\sigma}$. 

The kernel of (80.4) is denoted by

$$I_{\mathfrak{P}} = \{ \sigma \in G_{\mathfrak{P}} \mid \overline{\sigma} = 1_{A_{\mathfrak{P}}/\mathfrak{P}} \}$$

and called the inertia group of $\mathfrak{P}$, its fixed field

$$T_{\mathfrak{P}} = K^{I_{\mathfrak{P}}}$$

is called the inertia field of $\mathfrak{P}$. By Proposition 80.5, we have an isomorphism $G_{\mathfrak{P}}/I_{\mathfrak{P}} \cong G((A_{\mathfrak{P}})/(A/\mathfrak{p}))$. We know that $(A_{\mathfrak{P}})/(A/\mathfrak{p})$ is normal, but it is not Galois in general. However, if $K$ is a global field, it is as $A_{\mathfrak{P}}/\mathfrak{P}$ is then a finite field so separable over $(A/\mathfrak{p})$.

If the extension $(A_{\mathfrak{P}})/(A/\mathfrak{p})$ is separable, then $I_{\mathfrak{P}} = 1$ if and only if $\mathfrak{p}$ is unramified and $G_{\mathfrak{P}} \cong G((A_{\mathfrak{P}})/(A/\mathfrak{p}))$. More generally, in this case of separability, we have $e(\mathfrak{P}/\mathfrak{p}) = [I_{\mathfrak{P}}] = [K : T_{\mathfrak{P}}]$ and $f(\mathfrak{P}/\mathfrak{p}) = [G_{\mathfrak{P}} : I_{\mathfrak{P}}] = [T_{\mathfrak{P}} : \mathfrak{P}]$.

Suppose that $K$ is a global field, so $F = \mathbb{Q}$ and $A = \mathbb{Z}_K$ or $F$ a finite field and $A = F[t]_K$. Of course, if $F = \mathbb{Q}$, then $\mathfrak{p} = (p)$ in $\mathbb{Z}$ for some prime element $p$. Suppose that $\mathfrak{P}$ is unramified over $\mathfrak{p}$, then there exists a unique element $\sigma_{\mathfrak{P}}$ in $G(K/F)$ satisfying

$$\sigma_{\mathfrak{P}}(x) \equiv x^{f(\mathfrak{P}/\mathfrak{p})} \mod \mathfrak{P}.$$ 

The automorphism $\sigma_{\mathfrak{P}}$ is called the Frobenius automorphism. If $K/F$ is abelian, then $G_{(\mathfrak{P})} = \tau G_{\mathfrak{P}} \tau^{-1} = G_{\mathfrak{P}}$ for all $\tau$ in $G(K/F)$. As $G(K/F)$ acts transitively on the primes lying over $\mathfrak{p}$ all decomposition groups over $\mathfrak{p}$ are equal, so the Frobenius homomorphism $\sigma_{\mathfrak{P}}$ depends only on $\mathfrak{p}$. This was the first case, called ‘class field theory’, thoroughly studied in the first third of the twentieth century.

**Exercises 80.6.**

1. Let $A$ be a Dedekind domain with quotient field $F$, $\mathfrak{p}$ a nonzero prime ideal in $A$. Suppose that $K/F$ is a finite Galois extension and $\mathfrak{P}$ a
prime ideal in $A_K$ lying over $p$. Show that $T_{\mathfrak{p}}/Z_{\mathfrak{p}}$ is a normal extension satisfying $G(T_{\mathfrak{p}}/Z_{\mathfrak{p}}) \cong G(A_K/\mathfrak{p})/(A/p)$ and $G(K/T_{\mathfrak{p}}) \cong I_{\mathfrak{p}}$.

2. Let $A$ be a Dedekind domain with quotient field $F$, $p$ a nonzero prime ideal in $A$. Suppose that $K/F$ is a finite Galois extension and $\mathfrak{P}$ a prime ideal in $A_K$ lying over $p$ and satisfying $(A_K/\mathfrak{P})/(A/p)$ is separable. Assuming the previous exercise and Exercise 79.7(2) show

(i) $e(\mathfrak{P}/p) = [I_{\mathfrak{P}}] = [K : T_{\mathfrak{P}}]$ and $f(\mathfrak{P}/p) = [G_{\mathfrak{P}} : I_{\mathfrak{P}}] = [T_{\mathfrak{P}} : Z_{\mathfrak{P}}]$.
(ii) $e(\mathfrak{P}/\mathfrak{P} \cap T_{\mathfrak{P}}) = e(\mathfrak{P}/p)$ and $f(\mathfrak{P}/\mathfrak{P} \cap T_{\mathfrak{P}}) = 1$.
(iii) $e(\mathfrak{P} \cap T_{\mathfrak{P}}/\mathfrak{P} \cap Z_{\mathfrak{P}}) = 1$ and $f(\mathfrak{P} \cap T_{\mathfrak{P}}/\mathfrak{P} \cap Z_{\mathfrak{P}}) = f(\mathfrak{P}/p)$.

81. The Discriminant of a Number Field

In Section 63, we studied the discriminant of a polynomial $f$ in $F[t]$. If $K/F$ is a finite separable extension, then by the Primitive Element Theorem, $K = F(\alpha)$, for some $\alpha$. Of course there are many such $\alpha$’s, so there is no unique such discriminant. Suppose that $F$ is the quotient field of a Dedekind domain $A$. Then we know that we can choose $\alpha$ to lie in $A$. We now have two problems, $\alpha$ is still not unique, and more seriously, $A_K$ is not necessarily $A[\alpha]$ for any $\alpha$. We investigate what we must do to solve this problem, as the correct definition is crucial to finding the primes in $A$ that ramify in $A_K$. Although we shall not prove the full result, this addendum should provide a suitable introduction.

Let $K/F$ be a finite, separable field extension of degree $n$, $L/F$ a Galois extension with $L/K$, and $\sigma_1, \ldots, \sigma_n : K \to L$ the distinct $F$-homomorphisms. Let $w_1, \ldots, w_n$ be elements in $K$ and define

$$\Delta(w_1, \ldots, w_n) = \Delta_{K/F}(w_1, \ldots, w_n) := \det \left( \text{Tr}_{K/F}(w_iw_j) \right).$$

Then $\Delta(w_1, \ldots, w_n)$ lies in $F$ as each $\text{Tr}(w_iw_j)$ does. Moreover, if every $w_i$ lies in $A_K$, then $\Delta(w_1, \ldots, w_n)$ lies in $A$. We have (cf. also Section 63):

Properties 81.1. Let $K/F$ be a finite, separable field extension of degree $n$, $L/F$ a Galois extension with $L/K$, and $\sigma_1, \ldots, \sigma_n : K \to L$ the distinct $F$-homomorphisms. Let $B = \{w_1, \ldots, w_n\}$

(1) If $K = F(\alpha)$, then $\{1, \alpha, \ldots, \alpha^{n-1}\}$ is an $F$-basis for $K$. Set $\alpha_i = \sigma_i(\alpha)$. Then

$$\Delta(1, \alpha, \ldots, \alpha^{n-1}) = \prod_{i<j} (\alpha_i - \alpha_j)^2$$

is the discriminant of the minimal polynomial $m_F(\alpha)$. We denote it by $\Delta(\alpha)$. 


(2) \( \Delta(w_1, \ldots, w_n) \) is not zero if and only if \( B \) is an \( F \)-basis for \( K \).

(3) If \( C \) is an \( F \)-basis for \( K \) and \( \{w'_1, \ldots, w'_n\} \) its complementary basis, i.e., \( \text{Tr}_{K/F}(w_iw'_j) = \delta_{ij} \) for all \( 1 \leq i, j \leq n \), then
\[
\Delta(w_1, \ldots, w_n)\Delta(w'_1, \ldots, w'_n) = 1.
\]

(4) If \( C \) is an \( F \)-basis for \( K \), then
\[
\Delta(w_1, \ldots, w_n) = \det((\sigma_i(w_j))^2).
\]

(5) If \( B \) and \( \{v_1, \ldots, v_n\} \) are two \( F \)-bases for \( K \), then
\[
\Delta(v_1, \ldots, v_n) = (\det C)^2 \Delta(w_1, \ldots, w_n)
\]
with \( C \in \text{GL}_n(F) \), the change of basis matrix.

(6) If \( F = \mathbb{Q} \) and \( B \) and \( \{v_1, \ldots, v_n\} \) are two integral bases for \( \mathbb{Z}_K \), then
\[
\Delta(v_1, \ldots, v_n) = \Delta(w_1, \ldots, w_n),
\]
as the change of basis matrix for these two integral bases must lie in \( \text{GL}_n(\mathbb{Z}) \). Denote the integer \( \Delta(v_1, \ldots, v_n) \), independent of integral basis, by \( d_K \). It is called the discriminant of \( K \).

(7) Suppose that \( F = \mathbb{Q} \) and \( B \) is an integral basis for \( \mathbb{Z}_K \) and \( \{v_1, \ldots, v_n\} \) a \( \mathbb{Q} \)-basis for \( K \) with each \( v_i \) lying in \( \mathbb{Z}_K \). Then there exists an integer \( a \) such that \( a^2d_K = \Delta(v_1, \ldots, v_n) \) in \( \mathbb{Z} \); and, in fact,
\[
|\Delta(v_1, \ldots, v_n)| \geq |d_K|
\]
with equality if and only if
\[
\{v_1, \ldots, v_n\} \text{ is an integral basis.}
\]
[For this reason, an integral basis for \( \mathbb{Z}_K \) is often called a minimal basis.]

In the case of a ring of algebraic integers, we have the following:

**Proposition 81.2.** Let \( K/\mathbb{Q} \) be a finite field extension of degree \( n \) and \( \{\alpha_1, \ldots, \alpha_n\} \) a subset of \( \mathbb{Z}_K \) forming a \( \mathbb{Q} \)-basis for \( K \). Then we have
\[
\mathbb{Z}[\alpha_1, \ldots, \alpha_n] \subset \mathbb{Z}_K \subset \frac{1}{\Delta(\alpha_1, \ldots, \alpha_n)}\mathbb{Z}[\alpha_1, \ldots, \alpha_n].
\]

**Proof.** The first inclusion is trivial. Let \( \Delta = \Delta(\alpha_1, \ldots, \alpha_n) \) and \( \delta = \det((\sigma_i(\alpha_j)) \in \mathbb{Z}_K \), where \( \sigma_1, \ldots, \sigma_n : K \to \mathbb{C} \) are the \( \mathbb{Q} \)-embeddings of \( K \). So \( \Delta = \delta^2 \) by Property 81.1(4). Suppose that \( z \in \mathbb{Z}_K \). Then \( z = x_1\alpha_1 + \cdots + x_n\alpha_n \) for some \( x_1, \ldots, x_n \) in \( \mathbb{Q} \). We have
\[
\sigma_i(z) = x_1\sigma_i(\alpha_1) + \cdots + x_n\sigma_i(\alpha_n) \text{ for } i = 1, \ldots, n.
\]
By Cramer’s Rule, $x_i = y_i/\delta$ with $y_i$ the determinant of the matrix $(\sigma_i(\alpha_j))$ with the transpose $(\sigma_1(z), \ldots, \sigma_n(z))^t$, replacing its $i$th column for each $i$. Since $y_1, \ldots, y_n, \delta$ are algebraic integers, $\Delta x_i = \delta^2 x_i = \delta y_i$ lies in $\mathbb{Z}_K \cap \mathbb{Q} = \mathbb{Z}$, i.e., $x_i$ lies in $\frac{1}{\delta} \mathbb{Z}$ for $i = 1, \ldots, n$. The result follows. \hfill \Box

**Definition 81.3.** Let $A$ be a Dedekind domain with quotient field $F$, and $K/F$ a finite separable field extension. The **discriminant ideal** of $A_K$ is the ideal $\mathfrak{D}_{A_K/A}$ of $A$ generated by the elements $\Delta(w_1, \ldots, w_n)$ in $A$ for all $F$-bases $\{w_1, \ldots, w_n\}$ of $K$ lying in $A_K$. Of course, if $A$ is a ring of integers, then $\mathfrak{D}_{A_K/A} = (d_K)$, but in this case it is important to work with the integer $d_K$ as it has a unique sign attached to it.

**Proposition 81.4.** Let $A$ be a Dedekind domain with quotient field $F$ and $K/F$ a finite separable field extension. Suppose that $K = F(\alpha)$ with $\alpha$ an element of $A_K$. If $\mathfrak{p}$ is a nonzero prime ideal of $A$ relatively prime to both $\mathfrak{D}_{A_K/A}$ and the conductor $\mathfrak{f}$ of $A[\alpha]$ in $A_K$, then $\mathfrak{p}$ is unramified in $A_K$. In particular, only finitely many primes in $A$ ramify in $A_K$.

**Proof.** Let $\sim : A_K[t] \to (A_K/\mathfrak{p}A_K)[t]$ be the canonical epimorphism and

$$m_F(\alpha) = p_1^{e_1} \cdots p_r^{e_r}$$

be a factorization in $\overline{A}[t]$ into irreducible polynomials with $p_1, \ldots, p_r$ monic polynomials in $A[t]$. By Theorem 79.5,

$$\mathfrak{p}_i = \mathfrak{p}A_K + p_i(\alpha)A_K, \quad i = 1, \ldots, r$$

are the prime ideals lying over $\mathfrak{p}$ in $A_K$ and $pA_K = \mathfrak{p}_1^{e_1} \cdots \mathfrak{p}_r^{e_r}$. Let $\alpha_1, \ldots, \alpha_s$ be the roots of $m_F(\alpha)$ in a finite Galois extension of $F$ containing $K$. Then $\Delta(\alpha) = \prod_{i < j} (\alpha_i - \alpha_j)^2$, so $m_F(\alpha)$ has multiple roots if and only if $\alpha_i = \alpha_j$ for some $i \neq j$ if and only if $\Delta(\alpha) = 0$. Hence $e_i = 1$ for $i = 1, \ldots, r$ and each $\overline{p}_i$ has only simple roots. In particular, $A_K/\mathfrak{p}_i^{e_i} = A_K/\mathfrak{p}_i = \overline{A}_K(\overline{\alpha_i})$ is a separable extension of $\overline{A}$. \hfill \Box

Of course in the theorem, if $A_K = A[\alpha]$ for some $\alpha$, then $\mathfrak{f}$ is the unit ideal, so the condition in this case is that $\mathfrak{p}$ is relatively prime to $\mathfrak{D}_{A_K/A}$. Unfortunately, in general $A_K$ is not $A[\alpha]$ for any $\alpha$. However, in the number field case, we can show by elementary means that if a prime ideal $\mathfrak{p}\mathbb{Z}$ in $\mathbb{Z}$ ramifies in $\mathbb{Z}_K$ for $K/Q$ finite, then $\mathfrak{p}$ divides $d_K$.

**Proposition 81.5.** Let $K$ be a number field and $p$ a prime integer. If $(p)$ ramifies in $\mathbb{Z}_K$, then $p \mid d_K$. 
Proof. Let \( L/K \) be a finite extension with \( L/Q \) a finite Galois extension. Suppose that \( p \not| \, d_K \) and \( (p) \) ramifies in \( \mathbb{Z}_K \), say \( p\mathbb{Z}_K = \mathfrak{P}_1^{e_1} \cdots \mathfrak{P}_r^{e_r} \) is a factorization of \( p\mathbb{Z}_K \) with \( e_1 > 1 \). Let \( \mathfrak{A} = \mathfrak{P}_1^{e_1-1} \cdots \mathfrak{P}_r^{e_r} \) be a prime ideal in \( \mathbb{Z}_K \). We have \( \mathfrak{A}_i \mid \mathfrak{A} \) for \( i = 1, \ldots, r \). Let \( \{ x_1, \ldots, x_n \} \) be an integral basis for \( \mathbb{Z}_K \) where \( n = [K : Q] \). Choose \( x \in \mathfrak{A} \setminus p\mathbb{Z}_K \) and write

\[
x = m_1x_1 + \cdots + m_nx_n
\]

with \( m_1, \ldots, m_n \in \mathbb{Z} \). By assumption, there exists an \( i \) such that \( p \not| \, m_i \). We may assume that \( i = 1 \). We have \( \mathfrak{A} \cap \mathbb{Q} = \mathfrak{A} \). Choose \( x \) in \( \mathfrak{A} \setminus p\mathbb{Z}_K \) and write

\[
x = m_1x_1 + \cdots + m_nx_n
\]

As \( x \) lies in every prime ideal in \( \mathbb{Z}_K \) lying over \( (p) \), it lies in every prime ideal in \( \mathbb{Z}_L \) lying over \( (p) \). Let \( \mathfrak{Q} \) be a prime ideal in \( \mathbb{Z}_L \) lying over \( (p) \). As \( G(L/Q) \) acts transitively on the set of prime ideals in \( \mathbb{Z}_L \) lying over \( (p) \), we have \( \sigma(x) \) lies in \( \mathfrak{Q} \) for every automorphism \( \sigma \) in \( G(L/Q) \). By Property 81.1(4), the element \( \Delta(x, x_2, \ldots, x_n) \) lies in \( \mathfrak{Q} \cap \mathbb{Z} = (p) \). Consequently, \( p \mid \Delta(x, x_2, \ldots, x_n) = m_1^2\Delta(x_1, \ldots, x_n) \). Since \( p \not| \, m_1 \), we must have \( p \mid \Delta(x_1, \ldots, x_n) = d_K \), a contradiction. \( \square \)

It can be shown that the converse to the last proposition is in fact true, i.e., if \( p \) is a prime integer, then \( (p) \) ramifies in a ring of algebraic integers \( \mathbb{Z}_K \) if and only if \( p \mid d_K \). Indeed, in Section 82, we shall prove Dedekind’s Ramification Theorem that states if \( A \) is a Dedekind domain with quotient field \( F \) and \( K/F \) is a finite separable field extension, then a prime ideal \( \mathfrak{p} \) ramifies in \( A_K \) if and only if \( \mathfrak{p} \mid \mathfrak{D}_{A_K/F} \).

In the exercises for this section, we shall use localization techniques to reduce this to the case of a discrete valuation ring. In this section we shall, however, prove the following special case of it.

**Proposition 81.6.** Let \( p \) be an odd prime, \( \zeta \) a primitive \( p \)th root of unity in \( \mathbb{C} \), and \( K = \mathbb{Q}(\zeta) \). Then \( \mathbb{Z}_K = \mathbb{Z}[\zeta] \). In particular, \( \Delta_K = (-1)^{n-1}p^{p-2} \). Moreover, \( p\mathbb{Z} \) ramifies in \( \mathbb{Z}_K \) and is the only prime ideal that ramifies in \( \mathbb{Z}_K \).

**Proof.** By Example 63.7, we know that \( \Delta(\zeta) = (-1)^{n-1}p^{p-2} \) and that \( N(\zeta) = N(1 - \zeta) \). Certainly, \( \mathbb{Z}[\zeta] = \mathbb{Z}[1 - \zeta] \). Suppose that \( \mathbb{Z}[1 - \zeta] = \mathbb{Z}[\zeta] < \mathbb{Z}_K \). Using Proposition 81.2, we can find an element
z ∈ \mathbb{Z}_K \setminus \mathbb{Z}[1 - \zeta] satisfying
\[
z = \sum_{j=1}^{p-1} \frac{m_j}{p}(1 - \zeta)^j \quad \text{for some } i \geq 0 \text{ and } m_j \in \mathbb{Z} \text{ with } p \nmid m_i.
\]
(why?) By Example 63.7, we know that
\[
\prod_{j=0}^{p-1} (\zeta^j - 1) = p,
\]
hence \(p/(1 - \zeta)^{p-1}\) lies in \(\mathbb{Z}[\zeta]\), as \((1 - \zeta) \mid (1 - \zeta^j)\) for all \(j > 0\). Consequently, \(p/(1 - \zeta)^j\) lies in \(\mathbb{Z}[\zeta]\) for \(j = 0, \ldots, p - 1\) and \(zp/(1 - \zeta)^j\) lies in \(\mathbb{Z}_K\) for all \(j\). It follows easily that this implies that \(m_i/(1 - \zeta)\) lies in \(\mathbb{Z}_K\). But
\[
p = N_{K/Q}(1 - \zeta) \mid N_{K/Q}(m_i) = m_i^{p-1},
\]
which is impossible. Therefore, \(\mathbb{Z}_K = \mathbb{Z}[\zeta]\).

By Property 81.1(6), we have \(d_K = \Delta(\zeta)\). Therefore, \(\mathbb{Q}(\sqrt{\Delta}) = \mathbb{Q}\left(\sqrt{-\frac{p-1}{2}p}\right)\), the unique quadratic extension of \(\mathbb{Q}\) in \(K\) by Proposition 58.14. Clearly, \(p\mathbb{Z}\) ramifies in \(\mathbb{Z}_\mathbb{Q}(\sqrt{\Delta})\), so the proposition follows. \(\square\)

Remark 81.7. If \(\zeta\) is a primitive \(p^r\)th root of unity and \(\zeta_p\) a primitive \(p\)th root of unity in \(\mathbb{C}\) with \(p\) an odd prime or \(p = 2\) and \(r > 1\), then a modification of the proof shows that \(\mathbb{Z}_{\mathbb{Q}(\zeta)} = \mathbb{Z}[\zeta] = \mathbb{Z}[1 - \zeta]\) and the discriminant \(d_{\mathbb{Z}_{\mathbb{Q}(\zeta)}}\) is a power of \(p\) by Lemma 63.8. Using Proposition 81.5 and the fact that \(\mathbb{Q}\left(\sqrt{(-1)^{\frac{p-1}{2}p}}\right)\) is a subfield of \(\mathbb{Q}(\zeta)\) if \(p\) is odd with \(p\mathbb{Z}\) ramifying in \(\mathbb{Z}\left(\sqrt{\zeta^{p-1}}\right)\) and with \(2\) ramifying \(\mathbb{Z}[\sqrt{-1}]\) (cf. the next section), we see that such \(p\mathbb{Z}\) are precisely the prime ideals ramifying in \(\mathbb{Z}_K\). In the general case, one shows that \(\mathbb{Z}_L = \mathbb{Z}[\omega]\) for a primitive \(n\)th root of unity \(\omega\) for arbitrary \(n\). This is left as an exercise. The primes \(p\) that ramify in \(\mathbb{Z}_L\) are precisely those odd primes dividing \(n\) and \(2\) if \(4 \mid n\) using a similar analysis.

It can be shown that if \(\mathbb{Q} < K\) is a finite extension, then \(d_K > 1\) and some prime ideal in \(\mathbb{Z}\) ramifies in \(\mathbb{Z}_K\).

Exercises 81.8.
1. Prove Properties 81.1.
2. (Stickelberger’s Criterion) Let \(K\) be a number field. Show that \(d_K\) is congruent to 0 or 1 modulo 4.
3. (Brill) Let \(K\) be a number field and \(r\) the number of embedding of \(K\) into \(\mathbb{C}\) whose image does not lie in \(\mathbb{R}\). Show that the sign of \(d_K\) is \((-1)^r\).
4. Let \( \zeta \) be a primitive \( p^m \)th root of unity in \( \mathbb{C} \) with \( p \) an odd prime or \( p = 2 \) and \( m > 1 \). Let \( K = \mathbb{Q}(\zeta) \). Show that \( \Delta_K \mid p^{m-1}(p-1) \) and \( \mathbb{Z}_K = \mathbb{Z}[\zeta] = \mathbb{Z}[1 - \zeta] \). In particular, show if \( p \) is an odd prime or \( p = 2 \) and \( m \geq 2 \), then \( p\mathbb{Z} \) ramifies in \( \mathbb{Z}_K \) and is the only prime ideal ramifying in \( \mathbb{Z}_K \).

5. Let \( K/\mathbb{Q} \) and \( L/\mathbb{Q} \) be finite and \( d \) the gcd of \( d_K \) and \( d_L \). Show that \( \mathbb{Z}_K \mathbb{Z}_L \subset \frac{1}{d} \mathbb{Z}_{KL} \) where \( KL \) is the compositum \( K(L) - L(K) \) of \( K \) and \( L \) in \( \mathbb{C} \). In particular, if \( d = 1 \), then \( \mathbb{Z}_K \mathbb{Z}_L = \mathbb{Z}_{KL} \).

6. Using the last two exercises, show that if \( K = \mathbb{Q}(\zeta) \) with \( \zeta \) a primitive \( n \)th root of unity in \( \mathbb{C} \), then \( \mathbb{Z}_K = \mathbb{Z}[\zeta] \) and a prime \( p \) has \( p\mathbb{Z} \) ramifies in \( K \) if and only if \( p \mid n \) if \( p \) is odd or 2 if \( 4 \mid n \).

7. Let \( A \) be a Dedekind domain with quotient field \( F \) and \( K/F \) a finite separable field extension. Let \( \mathfrak{p} \) be a maximal ideal in \( A \) and \( S \) a multiplicative set not containing zero. Show that \( S^{-1}\mathfrak{D}_{A_K/A} = \mathfrak{D}_{S^{-1}A_K/S^{-1}A} \). In particular, if \( S = A \setminus \mathfrak{p} \), then \( \mathfrak{p} \mid \mathfrak{D}_{A_K/A} \) if and only if \( \mathfrak{p}S^{-1}A \mid \mathfrak{D}_{S^{-1}A_K/S^{-1}A} \). (Cf. Exercise 79.7(6.) [This reduces the theorem that \( \mathfrak{p} \) ramifies in \( A \) if and only if \( \mathfrak{p} \mid \mathfrak{D}_{A_K/A} \) if and only if it does so locally, i.e., \( A \) is a discrete valuation ring.]

8. Let \( A \) be a discrete valuation ring with quotient field \( F \) and \( K/F \) a finite separable field extension with \( A \)-basis \( \{x_1, \ldots, x_n\} \). Then \( \mathfrak{D}_{A_K/A} = \Delta_{K/F}(x_1, \ldots, x_n)A \).

**82. Dedekind’s Theorem on Ramification**

In this short section, we prove that a prime in a Dedekind domain ramifies in a finite separable extension of its quotient field if and only if it divides the discriminant ideal. We shall need to extend the definition of the trace of a finite separable field extension to an arbitrary finite field extension. We do this as follows: Let \( K/F \) be a finite field extension and \( K_s \) be the elements in \( K \) separable over \( F \). The set \( K_s \) is a field (why?) called the separable closure of \( F \) in \( K \). Define \( \text{Tr}_{K/F} : K \to F \) by \( \text{Tr}_{K/F} := [K : K_s] \text{Tr}_{K_s/F} \). In particular, \( K/F \) is not separable if and only if \( \text{Tr}_{K/F} = 0 \). (The norm \( N_{K/F} \) is defined by \( N_{K/F} := (N_{K_s/F})^{[K : K_s]} \).) Exercises in previous sections that entail the use of localization reduce the proof of our desired result to the case of a discrete valuation domain, i.e., a local Dedekind domain. In particular, assuming these exercises and this extension of the trace, we are reduced to proving the following result:

**Proposition 82.1.** Let \( A \) be a discrete valuation ring with quotient field \( F \) and \( K/F \) a finite separable extension. Then the maximal ideal \( \mathfrak{p} \) in \( A \) ramifies in \( K \) if and only if \( \mathfrak{p} \mid \mathfrak{D}_{A_K/F} \).
**Proof.** Let \( n = [K : F] \). Since a discrete valuation ring is a PID by Remark 78.8(5), we know that \( A_K \) is a free \( \mathcal{A} \)-module of rank \( n \), say with integral basis \( \{x_1, \ldots, x_n\} \). It follows by Exercise 81.8(8) that \( \mathcal{D}_{A_K/A} = \Delta_{K/F}(x_1, \ldots, x_n)A \). Hence \( p \mid \mathcal{D}_{A_K/A} \) if and only if \( \Delta_{K/F}(x_1, \ldots, x_n) \in p \).

Let \((a)\) \( \overline{A} = A/p \).

\((b)\) \( pA_K = \mathfrak{p}_1^{e_1} \cdots \mathfrak{p}_r^{e_r} \) be the factorization of \( p \) in \( A_K \) into primes.

\((d)\) \( B_i = A_K/\mathfrak{p}_i^{e_i} \) for \( i = 1, \ldots, r \).

Using the Chinese Remainder Theorem, we identify \( A_K/pA_k \) with \( \prod_{i=1}^{r} B_i \).

We know that \( n = [K : F] = \text{rank}_A(A_K) = \dim_{\overline{A}} \overline{A_K} \)

by Theorem 79.2 (and its proof). Let \( \mathcal{B}_i = \{u_{j_i-1+1}, \ldots, u_{j_i}\} \) be an \( \overline{A} \)-basis for \( B_i \), \( i = 1, \ldots, r \). So \( \mathcal{B} = \mathcal{B}_1 \cup \cdots \cup \mathcal{B}_r \) is an \( \overline{A} \)-basis for \( \overline{A_K} \) with \( |\mathcal{B}| = n \).

For each \( y \) in \( \overline{A_K} \) define \( \lambda_y : \overline{A_K} \to \overline{A_K} \) by \( z \mapsto yz \) and \( \text{tr} : \overline{A_K} \to \overline{A_K} \) by \( y \mapsto \text{tr}(y) := \text{trace}(\lambda_y) \). Let \( \text{tr}_i = \text{tr}|_{B_i} \) for \( i = 1, \ldots, r \).

Note that we have

\((i)\) \( \{\overline{x}_1, \ldots, \overline{x}_n\} \) is an \( \overline{A} \)-basis for \( \overline{A_K} \).

\((ii)\) \( \text{tr}(u_iu_j) = 0 \) if \( u_i \) and \( u_j \) lie in different \( B_k \)'s.

**Claim.** Let \( x \in A_K \), then \( \overline{\text{Tr}}_{K/F}(x) = \text{tr}(\overline{x}) \) \( (= \text{trace}(\lambda_{\overline{x}})) \). In particular, there exists a change of basis matrix \( C \in \text{GL}_n(\overline{A}) \) satisfying

\[ \Delta_{K/F}(x_1, \ldots, x_n) = (\det(C))^2 \det(\text{tr}(u_iu_j)), \]

hence

\[ \Delta_{K/F}(x_1, \ldots, x_n) = 0, \text{ i.e., } \Delta_{K/F}(x_1, \ldots, x_n) \in p, \]

if and only if \( \det(\text{tr}(u_iu_j)) = 0 \):

Write \( xx_i = \sum_{j=1}^{n} a_{ij}x_j \) with \( a_{ij} \in A \) for all \( i \) and \( j \) and let \( \alpha \) be the matrix \( (a_{ij}) \). Then \( \alpha \) is the matrix representation of \( \lambda_x \) relative to the basis \( \{x_1, \ldots, x_n\} \) and \( \overline{\alpha} \) is the matrix representation of \( \lambda_{\overline{x}} \) relative to the basis \( \{\overline{x}_1, \ldots, \overline{x}_n\} \). By Exercise 59.26(5), we have

\[ \overline{\text{Tr}}_{K/F}(x) = \text{trace}(\lambda_{\overline{x}}) = \text{trace}(\lambda_{\overline{x}}) = \text{tr}(\overline{x}). \]
This proves the Claim. Now let \( D_i := (\text{tr}_i(u_{li}u_{lk})) \) for \( j_{i-1} < l, k \leq j_i \), for each \( i = 1, \ldots, r \). Then

\[
D := (\text{tr}_i(u_{li}u_{lk})) = \begin{pmatrix}
D_1 & \cdots & 0 \\
\vdots & \ddots & \vdots \\
0 & \cdots & D_r
\end{pmatrix} \quad \text{(with } j_0 = 0).\]

Suppose that \( B_i = A_K / \mathfrak{P}_i^{e_i} \) is a field, i.e., \( e_i = 1 \). Then

\[
det D_i = \Delta_{B_i/\mathfrak{A}}(u_{j_{i-1}+1}, \ldots, u_{j_i}) = \det \left( \text{Tr}_{B_i/\mathfrak{A}}(u_{li}u_{lk}) \right).
\]

Since \( B_i/\mathfrak{A} \) is not separable if and only if \( \text{Tr}_{B_i/\mathfrak{A}} = 0 \), if follows that if \( B_i \) is a field for every \( i \), then \( p | D_AK/A \) if and only \( \det D_i = 0 \) for some \( i \) if and only if there exists an \( i \) such that \( B_i/\mathfrak{A} \) is not separable.

So suppose that there exists an \( i \) with \( B_i \) not a field. We may assume that \( i = 1 \) and further that \( u_1 \) lies in \( \mathfrak{P}_1/\mathfrak{P}_1^{e_1} \). Therefore, \( (u_1u_i)^{e_1} = 0 \) for all \( i \), hence \( \lambda_{u_1u_i} \) is nilpotent. It follows that the first row of \( D \) consists of 0’s by Example 44.7, hence \( \det D = 0 \) and \( p | D_AK/A \).

The generalization of the proposition to arbitrary Dedekind domains now follows by the reduction given by Exercise 81.8(7). That is, we the following theorem:

**Theorem 82.2.** (Dedekind Ramification Theorem) Let \( A \) be a Dedekind domain with quotient field \( F \) and \( K/F \) a finite separable field extension. The a maximal ideal \( p \) of \( A \) ramifies in \( A_K \) if and only if \( p | D_AK/A \).

Immediate consequences of the theorem are:

**Corollary 82.3.** Let \( A \) be a Dedekind domain with quotient field \( F \) and \( K/F \) a finite separable field extension. Then \( A_K/A \) is unramified if and only if \( D_AK/A = A \).

**Remarks 82.4.** Let \( F \) be a number field, so a finite extension of \( \mathbb{Q} \).

1. Using Minkowski Theory, one can show that \( |d_F| > 1 \), so \( \mathbb{Z}_F/\mathbb{Z} \) cannot be unramified. As \( \mathbb{Q} \) and finite fields are perfect, this means that there exists a prime \( p \) and a prime ideal \( \mathfrak{P} \) in \( \mathbb{Z}_F \) lying over \( p\mathbb{Z} \) such that \( e(\mathfrak{P}/p\mathbb{Z}) > 1 \).

2. It can also be shown that there exist finitely many \( F/\mathbb{Q} \) satisfying \( d = d_F \), for some fixed integer \( d \).

3. There can, however, be finite extensions \( K \) of the number field \( F \) that are unramified. Indeed a deep theorem of Classfield Theory says that the maximal unramified abelian extension \( K \) of \( F \) (in \( \mathbb{C} \)) is finite and of degree \( |Cl_{\mathbb{Z}_K}| \) over \( F \), where \( Cl_{\mathbb{Z}_F} (= I_{\mathbb{Z}_F}/P_{\mathbb{Z}_K}) \) is the class group of \( \mathbb{Z}_K \) (finite for a ring of algebraic integers by Minkowski...
Theory – cf. Section ??). Furthermore, its Galois group $G(K/F)$ is (canonically) isomorphic to $CL\mathbb{Z}_F$. [K is called the Hilbert class field of F.] In particular, $F < K$ if $\mathbb{Z}_F$ is not a PID.

**Corollary 82.5.** Let $F = \mathbb{Q}(\alpha)$ with $\alpha$ integral over $\mathbb{Z}$ satisfying $g(\alpha) = 0$ with $g \in \mathbb{Z}[t]$ monic and $p$ a prime in $\mathbb{Z}$. If $p \nmid N_{F/\mathbb{Q}}(g'(\alpha))$, then $p\mathbb{Z}$ is unramified in $\mathbb{Z}_K$.

**Proof.** Let $f = m_{\mathbb{Q}}(\alpha)$. By Proposition 63.6, we know that $\Delta_{F/\mathbb{Q}}(\alpha) = \pm N_{\mathbb{F}/\mathbb{Q}}(f'(\alpha))$. Let $g = fh$ in $\mathbb{Z}[t]$, then $g' = f'h + fh'$. Consequently, $d_F \mid \Delta_{\mathbb{F}/\mathbb{Q}}(\alpha) N_{\mathbb{F}/\mathbb{Q}}(h(\alpha)) = \pm N_{\mathbb{F}/\mathbb{Q}}(g'(\alpha))$.

### 83. The Quadratic Case

Let $d$ be a square-free integer. In this section, we compute the ring theory of $\mathbb{Z}_K$ with $K = \mathbb{Z}[\sqrt{d}]$ to illustrate the theory that we have developed.

Throughout this section let $\mathbb{Z}_K$ with $K = \mathbb{Z}[\sqrt{d}]$ and $G(\mathbb{K}/\mathbb{F}) = \{1, \sigma\}$, the Galois group of the Galois extension $K/\mathbb{Q}$. We first determine an integral basis and discriminant of $K$. We freely use Properties 81.1. We need a special case of Stickelberger’s Criterion (Exercise 81.8(2)), viz.,

$$d_K \equiv 0 \mod 4 \text{ or } d_K \equiv 1 \mod 4.$$ 

Indeed, if $\{w_1, w_2\}$ is an integral basis for $\mathbb{Z}_K$, let $\alpha = w_1\sigma(w_2)$ and $\beta = w_2\sigma(w_1)$. Then

$$d_K = (\alpha - \beta)^2 = (\alpha + \beta)^2 - 4\alpha\beta.$$ 

As $(\alpha - \beta)^2$, $(\alpha + \beta)^2$ and $\alpha\beta$ are all fixed by $G(\mathbb{K}/\mathbb{Q})$, they are all integers, and the result follows. (The proof in the general case is analogous.) The above notation shall be used throughout this section. We also have:

$$\Delta(1, \sqrt{d}) = \det \begin{pmatrix} 1 & \sqrt{d} \\ 1 & \sigma(\sqrt{d}) \end{pmatrix}^2 = \det \begin{pmatrix} 1 & \sqrt{d} \\ 1 & -\sqrt{d} \end{pmatrix}^2 = 4d.$$

As $\mathbb{Z}[\sqrt{d}] \subset \mathbb{Z}_K$, there exists an integer $a \in \mathbb{Z}$, satisfying

$$a^2d_K = 4d, \text{ so } d_K = 4d \text{ or } d_K = d,$$

since $d$ is square-free. (Note if $a = 2$, we must have $d$ is odd, lest $d_K \equiv 2 \mod 4$ which contradicts Stickelberger’s Criterion. So if $d$ is even, $d_K = 4d$ and $8 \mid d_K$.) Using this and Stickelberger’s Criterion, we solve our first goal.

**Case 1.** $d \equiv 2$ or $3 \mod 4$: 

We have \( d_K = 4d \) and \( \mathbb{Z}_K = \mathbb{Z}[\sqrt{d}] \). In particular, \( \{1, \sqrt{d}\} \) is an integral basis for \( \mathbb{Z}_K \) by Property 81.1(7), and the conductor \( f \) of \( \mathbb{Z}[\sqrt{d}] \) in \( \mathbb{Z}_K \) is \( \mathbb{Z}_K \).

**Case 2.** \( d \equiv 1 \mod 4 \): We have \( (1 - d)/4 \) is an integer, hence \( (1 \pm \sqrt{d})/2 \) lies in \( \mathbb{Z}_K \), as they are the roots of \( t^2 - t + (1 - d)/4 \) in \( \mathbb{Z}[t] \). It follows that \( \mathbb{Z}[\sqrt{d}] \subset \mathbb{Z}[\frac{1+\sqrt{d}}{2}] \subset \mathbb{Z}_K \).

We know that \( \Delta(1, (1 + \sqrt{d})/2) = b^2d_K \), some integer \( b \) and computation yields \( \Delta(1, (1 + \sqrt{d})/2) = d \). Therefore, by Property 81.1(7), \( d_K = d \), \( \mathbb{Z}_K = \mathbb{Z}[(1 + \sqrt{d})/2] \), and \( \{1, (1 + \sqrt{d})/2\} \) is an integral basis for \( \mathbb{Z}_K \). In particular, the conductor \( f \) of \( \mathbb{Z}[\sqrt{d}] \) in \( \mathbb{Z}_K \) contains the ideal \( (2, 1 - \sqrt{d}) \). We show \( f = (2, 1 - \sqrt{d}) \). The element \( (1 + \sqrt{d})/2 \) does not lie \( \mathbb{Z}[\sqrt{d}] \) hence does not lie in \( f \), since \( f \subset \mathbb{Z}[\sqrt{d}] \). It follows if \( a+b((1+\sqrt{d})/2) \) lies in \( f \) with \( a, b \in \mathbb{Z} \), we may assume that \( a = b = 1 \). As \( 1+(1+\sqrt{d})/2 \) does not lie in \( \mathbb{Z}[\sqrt{d}] \), we have \( f = (2, 1 - \sqrt{d}) \). If a prime ideal \( \mathfrak{p} \mid (1 - \sqrt{d}) \), then \( \mathfrak{p} \mathfrak{p}^2 \mid (1 - d)\mathbb{Z} \). It follows that no factor of \( p\mathbb{Z}_K \) with \( p \) an odd prime can be a factor of \( f \).

Because of this computation, if \( p \) is a prime integer we can compute the splitting behavior of \( p\mathbb{Z}_K \) in \( \mathbb{Z}_K \) for all primes \( p \) by Theorem 79.5 except if \( p = 2 \) and \( d \equiv 1 \mod 4 \). In particular, it is applicable to all odd primes.

**Case of an odd prime.** Let \( p \) be an odd prime. We know that \( 2 = [K : \mathbb{Q}] = efr \) with \( e \) the ramification index of primes over \( (p) \), \( f \) the inertia index of \( (p) \), and \( r \) the number of primes in \( \mathbb{Z}_K \) lying over \( (p) \). We also know that the Legendre symbol \( (\frac{d}{p}) \) determines whether \( t^2 - d \) splits over \( \mathbb{Z}/p\mathbb{Z}[t] \) or not when \( d \) is not divisible by \( p \), i.e., determines \( f \) for such \( p \). Using the Kummer-Dedekind Theorem 79.5 and Proposition 81.4, we conclude, if \( p \) is an odd prime, that

1. If \( p \mid d_K \), then \( t^2 - d = t^2 \) in \( \mathbb{Z}/p\mathbb{Z}[t] \), so
   \[ p\mathbb{Z}_K = (p, \sqrt{d})^2. \]
   Therefore \( (p) \) ramifies in \( \mathbb{Z}_K \) if and only if \( p \mid d_K \).
2. The ideal \( p\mathbb{Z}_K \) splits completely if and only if \( (\frac{d}{p}) = 1 \). If this is the case, then
   \[ p\mathbb{Z}_K = (p, x + \sqrt{d})(p, x - \sqrt{d}) \]
a product of two distinct prime ideals with the integer \( x \) satisfying 
\[ x^2 \equiv d \mod p. \]

3. The ideal \( p\mathbb{Z}_K \) remains prime in \( \mathbb{Z}_K \) if and only if \( (\frac{d}{p}) = -1. \)

**Case of the even prime 2:**

If \( d \equiv 2 \) or \( 3 \mod 4 \), then Theorem 79.5 applies, and we see that in either case (2) ramifies (note \( 2 \mid d_K \)), because

1. If \( d \equiv 2 \mod 4 \), then \( 2\mathbb{Z}_K = (2, \sqrt{d})^2. \)
2. If \( d \equiv 3 \mod 4 \), then \( 2\mathbb{Z}_K = (2, d + \sqrt{d})^2 = (2, 1 + \sqrt{d})^2. \)

So we are reduced to the case that \( d \equiv 1 \mod 4 \), and this reduces to the two cases of \( d \equiv 1 \mod 8 \) and \( d \equiv 5 \mod 8 \). Since Theorem 79.5 does not apply, we do this by brute force. Notice the computation shows that 2 does not ramify as expected and \((1 + \sqrt{d})/2\) lies in \( \mathbb{Z}_K \).

We have

3. If \( d \equiv 1 \mod 8 \), then (2) splits completely in \( \mathbb{Z}_K \):

   We have
   \[
   (2, \frac{1 + \sqrt{d}}{2})(2, \frac{1 - \sqrt{d}}{2}) = (4, 1 + \sqrt{d}, 1 - \sqrt{d}, \frac{1 - d}{4}) \subset 2\mathbb{Z}_K
   \]

   As \( 2\mathbb{Z}_K \subset (1 + \sqrt{d}, 1 - \sqrt{d}) \), we conclude that
   \[
   2\mathbb{Z}_K = (2, \frac{1 + \sqrt{d}}{2})(2, \frac{1 - \sqrt{d}}{2})
   \]
   splits completely.

4. If \( d \equiv 5 \mod 8 \), then (2) remains a prime ideal in \( \mathbb{Z}_K \):

   We must show that \( \mathbb{Z}/2\mathbb{Z} < \mathbb{Z}_K / \mathfrak{p} \) if \( \mathfrak{p} \mid 2 \). If this was false, there are only two cosets of \( \mathbb{Z}_K / \mathfrak{p} \), so there exists an integer \( a \) satisfying \( a \equiv (1 + \sqrt{d})/2 \mod \mathfrak{p} \). Since \((1 + \sqrt{d})/2\) is a root of \( t^2 - t + (1 - d)/4 \), we have \( a^2 - a + (1 - d)/4 \) is even. As \( a^2 - a \) is even, we would have \( d \equiv 1 \mod 8 \), a contradiction. Therefore, (2) remains a prime ideal.

As previously mentioned, a theorem of Kronecker says that any abelian extension of \( \mathbb{Q} \) lies in a cyclotomic extension of \( \mathbb{Q} \). We have shown that every quadratic extension of \( \mathbb{Q} \) is a subfield of a cyclotomic extension of \( \mathbb{Q} \) in Theorem 58.17. We use the proposition that implied this together with Hilbert Ramification Theory to interpret when an ideal \((p)\) with \( p \) an odd prime in \( \mathbb{Z} \) splits completely in \( \mathbb{Z}_K \).

**Proposition 83.1.** Let \( p \) and \( l \) be odd primes, \( L = \mathbb{Q}(\zeta) \) with \( \zeta \) a primitive \( l \) th root of unity, and \( K = \mathbb{Q}(\sqrt{\frac{-1}{l}}) \). Then \((p)\) splits
completely in \( \mathbb{Z}_K \) if and only if \((p)\) splits into an even number of prime ideals in \( \mathbb{Z}_L \).  

**Proof.** By Proposition 58.14, we have \( K \) is the unique quadratic extension of \( \mathbb{Q} \) in \( L \). Suppose that \( p \) splits completely in \( \mathbb{Z}_K \), say \( p\mathbb{Z}_K = \mathfrak{P}_1\mathfrak{P}_2 \) is its factorization. There exists an automorphism \( \sigma \in G(L/\mathbb{Q}) \) satisfying \( \sigma(\mathfrak{P}_1) = \mathfrak{P}_2 \) by Proposition 80.1. Then \( \sigma \) takes the set of prime ideals in \( \mathbb{Z}_L \) over \( \mathfrak{P}_1 \) bijectively onto the set of prime ideals over \( \mathfrak{P}_2 \). Therefore the set of primes over \((p)\) in \( \mathbb{Z}_L \) is even. Conversely, suppose that the number of prime ideals in \( \mathbb{Z}_L \) over \((p)\) is even. Let \( \mathfrak{P} \) be one such. Then the decomposition group \( G_\mathfrak{P} \) of \( \mathfrak{P} \) in \( G(L/\mathbb{Q}) \) has even index in \( G(L/\mathbb{Q}) \). By Galois Theory, this means the degree of the decomposition field \( \mathbb{Z}_\mathfrak{P} \) over \( \mathbb{Q} \) is even. But \( L/\mathbb{Q} \) is cyclic and even, so \( K \subset \mathbb{Z}_\mathfrak{P} \) using Proposition 80.3(2). \( \square \)  

Using the fact that if \( \zeta \) is a primitive \( n \)th root of unity that \( \mathbb{Z}_{\mathbb{Q}(\zeta)} = \mathbb{Z}[\zeta] \), a fact that we did not prove but left as exercises (cf. Exercise 81.8(6)), one can give another proof of the Law of Quadratic Reciprocity essentially due to Kronecker. Although we do not do this, it is key to the generalization of the Law of Quadratic Reciprocity in Number Theory called Artin Reciprocity. [Cf. Ireland and Rosen, *A Classical Introduction to Modern Number Theory* [20], Chapter 13, \( \S 3 \) for this proof.]  

We have determined the ideal structure of the quadratic number field \( \mathbb{Z}_K \) and from our work how nonzero ideals in \( \mathbb{Z} \) decompose in \( \mathbb{Z}_K \) as \( \mathbb{Z}_K \) is a Dedekind domain and ideals are products of prime ideals, unique up to order. The remaining problem is to determine the units in \( \mathbb{Z}_K \). This depends on whether \( d \) is positive or negative. The case of negative \( d \) is easy, the case of positive \( d \) not so.  

**Case of \( d < 0 \).** We show that  
1. If \( d = -1 \), then \( \mathbb{Z}_K^\times = \{\pm 1, \pm \sqrt{-1}\} \).  
2. If \( d = -3 \), then \( \mathbb{Z}_K^\times = \{\pm 1, \pm \zeta, \pm \zeta^2\} \), with \( \zeta \) a primitive cube root of unity.  
3. If \( d < -3 \) or \( d = -2 \), then \( \mathbb{Z}_K^\times = \{\pm 1\} \).  

We have seen this for the Gaussian integers, and the proof generalizes, i.e., we take the norm for \( K \) to \( \mathbb{Q} \). Let \( \varepsilon \) be a unit in \( \mathbb{Z}_K \). Then \( N_{\mathbb{K}/\mathbb{F}}(\varepsilon) = \pm 1 \).  

If \( d \equiv 2 \) or 3 modulo 4, we can write \( \varepsilon = x + y\sqrt{d} \) for some integers \( x \) and \( y \), so \( x^2 + |d|y^2 = 1 \). If \( d = -1 \), we get (1) and if \( |d| > 3 \) we have \( \mathbb{Z}_K^\times = \{\pm 1\} \).
If \( d \equiv 1 \mod 4 \), then \( \{1, (1 + \sqrt{d}/2)\} \) is an integral basis, and we see that we can write \( \varepsilon \) as \((x + y\sqrt{d})/2 \) with \( x \) and \( y \) integers satisfying \( x \equiv y \mod 2 \). We then have \( x^2 + |d|y^2 = 4 \). If \( d = -3 \) we get (2) and if \( |d| > 3 \), we see that \( \mathbb{Z}_K^\times = \{\pm 1\} \).

To do the case for positive \( d \), we need two results. The first is the following:

**Observation 83.2.** If \( M \) is a positive real number, there exist finitely many \( \alpha \) in \( \mathbb{Z}_K \) such that \( \max(\alpha, |\sigma(\alpha)|) < M \).

and the second a famous theorem from elementary number theory, viz.,

**Proposition 83.3.** Let \( d \) be a square-free positive integer. Then Pell’s equation \( x^2 - dy^2 = 1 \) has infinitely many solutions in integers. Further, there exists an integral solution \((x_1, y_1)\) such that every solution is of the form \((x_n, y_n)\) with \( x_n + y_n\sqrt{d} = (x_1 + y_1\sqrt{d})^n \) for some integer \( n \). In particular, the solution set in integers to Pell’s equation forms an infinite cyclic group.

that we prove in Appendix F.

Using these two results, we prove the case of positive discriminant.

**Case of \( d > 1 \).** We show that there exists a unit \( u > 1 \) in \( \mathbb{Z}_K \) such that \( (\mathbb{Z}_K)^\times = \{\pm u^m \mid m \in \mathbb{Z}\} \).

By Proposition 83.3, there exist positive integers \( x \) and \( y \) satisfying \( x^2 - dy^2 = 1 \). Let \( \varepsilon = x + y\sqrt{d} \). Then \( \varepsilon \) is a unit in \( \mathbb{Z}_K \) with \( \varepsilon > 1 \). Let \( M > \varepsilon \) and \( u \) be a unit in \( \mathbb{Z}_K \) satisfying \( 1 < u < M \). So \( N_{K/Q} = u\sigma(u) = \pm 1 \). If \( \sigma(u) = -1/u \), then \( -M < -1/u < M \) and if \( \sigma(u) = 1/u \), then we also have \( -M < 1/u < M \). By the observation, there exist only finitely many such \( u \) and \( \varepsilon \) is at least one of them. Among all of these choose the smallest unit \( u > 1 \). If \( v > 0 \) is any other unit, then there exist a unique integer \( n \) (not necessarily positive) satisfying \( u^n \leq v < u^{n+1} \), so \( 1 \leq vu^{-n} < u \). By the choice of \( u \), we must have \( v = u^n \). If \( v < 0 \) is a unit in \( \mathbb{Z}_K \), then \(-v > 0 \) is a unit in \( \mathbb{Z}_K \) and \(-v = u^n \) some integer \( n \). The result follows. [Note the similarity of the proof with that in Appendix F.]

The generator of the unit group in this case is called the *fundamental unit*. It can be found using continued fractions, but may not be so obvious. For example, the fundamental unit of \( \mathbb{Z}[\sqrt{94}] \) is 2143295 + 221064\sqrt{94}.

If \( K \) is a number field, then the structure of \( \mathbb{Z}_K \) can be shown to be \( \mu_K \times \mathbb{Z}^n \) where \( n = r + s - 1 \), with \( r \) the number of embeddings of \( K \) into \( \mathbb{R} \) and \( 2s \) the number of embeddings of \( K \) into \( \mathbb{C} \) not contained.
in \( \mathbb{R} \), so the number of embeddings into \( \mathbb{C} \) is \( r + 2s \) (as the non-real embeddings come in complex conjugate pairs). This is proved using Minkowski Theory in Section ??.

**Exercise 83.4.** Prove Observation 83.2.

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### 84. Addendum: Valuation Rings and Prüfer Domains

In this addendum, we shall generalize the concept of Dedekind domain. By Exercise 78.14(14), a Dedekind domain is a Noetherian domain \( A \) in which the localization at every nonzero prime ideal \( p \) in \( A \) is a discrete valuation ring. Another characterization of a Dedekind domain was that it was a Noetherian domain \( A \) in which every fractional ideal in \( A \) was invertible. We look at these when we do not assume the domain is Noetherian.

In Exercise 78.14(13) we defined a valuation ring. We begin by recalling the definition.

**Definition 84.1.** A domain \( A \) with quotient field \( F \) is called a valuation ring of \( F \) if for every nonzero element \( x \) in \( F \), either \( x \) lies in \( A \) or \( x^{-1} \) lies in \( A \). In particular, if \( a \) and \( b \) are nonzero elements in \( A \), then \( a \mid b \) or \( b \mid a \) in \( A \). We call a domain a valuation ring if it is a valuation ring in its quotient field.

Every discrete valuation ring is a Noetherian valuation ring by Exercise 78.14(13), hence the localization of any Dedekind domain at a nonzero prime ideal is a valuation ring. Valuation rings are ubiquitous, but unfortunately a valuation ring is Noetherian only when it is a discrete valuation ring. The general concept, however, was very important in the algebraization of algebraic geometry.

**Properties 84.2.** Let \( A \) be a valuation ring of \( F \). Then

1. \( A \) is a local ring.
2. If \( a_1, \ldots, a_n \) are nonzero elements in \( A \), then there exists an \( i \) such that \( a_i \mid a_j \) for \( j = 1, \ldots, n \).
3. If \( A \subset B \subset F \) are subrings, then \( B \) is a valuation ring.
4. \( A \) is integrally closed.
5. Every finitely generated ideal in \( A \) is principal. A domain in which every finitely generated ideal is principal is called a Bézout domain.

**Proof.** (1): Let \( \mathfrak{a} = A \setminus A^\times \), the set of nonunits in \( A \). As \( A \) is a valuation ring, \( a \) lies in \( \mathfrak{a} \) if and only if \( a^{-1} \) does not lie in \( \mathfrak{a} \).

**Claim.** \( \mathfrak{a} \) is an ideal.
Let \( a, b \) be nonzero elements in \( \mathfrak{A} \) and \( r \) a nonzero element in \( \mathfrak{A} \). Certainly, \( ra \in \mathfrak{A} \). Suppose that \( ra \notin \mathfrak{A} \). Then \( (ra)^{-1} \in \mathfrak{A} \), hence \( a^{-1} = r(ra)^{-1} \) lies in \( \mathfrak{A} \), a contradiction. So \( \mathfrak{A} \) is closed under multiplication by elements in \( \mathfrak{A} \). As either \( ab^{-1} \) or \( a^{-1}b \) lies in \( \mathfrak{A} \), say \( ab^{-1} \), it follows that \( a + b = b(ab^{-1} + 1) \) lies in \( \mathfrak{A} \). Hence \( \mathfrak{A} \) is an ideal and clearly the unique maximal ideal.

Statements (2) and (3) are immediate, (4) follows by Exercise 84.27(1), and (5) follows form (2). \( \square \)

We conclude the following:

**Corollary 84.3.** Let \( \mathfrak{A} \) be a Noetherian domain, not a field. Then \( \mathfrak{A} \) is a valuation ring if and only if \( \mathfrak{A} \) is a discrete valuation ring.

**Proof.** (\( \Rightarrow \)): By the Properties 84.2, \( \mathfrak{A} \) is a local PID so every nonzero prime ideal is maximal and integrally closed. Since it is Noetherian, it is a local Dedekind domain, hence a discrete valuation ring.

(\( \Leftarrow \)): Is Exercise 78.14(13). A key to the proof is if \( \mathfrak{m} \) is the maximal ideal of \( \mathfrak{A} \), then \( x \notin A^\times \cup \mathfrak{m}^2 \) must satisfy \( (x) = \mathfrak{m} \) by uniqueness of factorization of ideals into prime ideals in a Dedekind domain. It follows that \( \mathfrak{A} \) must be a PID and the result follows. \( \square \)

We next look at fractional ideals in a domain a bit more closely.

Let \( \mathfrak{A} \) be a domain with quotient field \( F \) and \( M \) a fractional ideal, i.e., a nonzero \( \mathfrak{A} \)-submodule of \( F \) such that there exists a nonzero element \( x \) of \( \mathfrak{A} \) satisfying \( xM \subset \mathfrak{A} \). As before, we shall let \( I_{\mathfrak{A}} \) denote the set of fractional ideals in \( \mathfrak{A} \). We say a fractional ideal \( M \) is an *invertible fractional ideal* if there exists an \( \mathfrak{A} \)-submodule \( N \) of \( F \) such that \( MN = \mathfrak{A} \). Let \( Inv(\mathfrak{A}) \) denote the set of invertible fractional ideals. As before, we let \( P_\mathfrak{A} = \{ Ax \mid 0 \neq x \in F \} \), the set of principal fractional ideals.

**Remarks 84.4.** Let \( \mathfrak{A} \) be a domain with quotient field \( F \) and \( M \) an \( \mathfrak{A} \)-submodule of \( F \).

1. Every nonzero ideal of \( \mathfrak{A} \) is a fractional ideal.
2. A nonzero finitely generated \( \mathfrak{A} \)-submodule of \( F \) is a fractional ideal — clear denominators.
3. Let \( x \) be a nonzero element of \( F \), then \( M \in Inv(\mathfrak{A}) \) if and only if \( xM \in Inv(\mathfrak{A}) \). In particular, \( P_\mathfrak{A} \subset Inv(\mathfrak{A}) \).
4. If \( \mathfrak{A} \) is Noetherian and \( M \in I_\mathfrak{A} \), then there exists a nonzero \( x \) in \( F \) such that \( xM \) is an ideal. In particular, \( M \) is a finitely generated \( \mathfrak{A} \)-module. It follows that if \( \mathfrak{A} \) is Noetherian, \( I_\mathfrak{A} \) is the set of nonzero finitely generated \( \mathfrak{A} \)-submodules of \( F \).
As before if \( M \in I_A \), we let \( M^{-1} := \{ x \in F^\times \mid xM \subset A \} \). Of course, \( MM^{-1} \subset A \).

**Lemma 84.5.** Let \( A \) be a domain with quotient field \( F \) and \( M, N \) be \( A \)-submodules of \( F \) satisfying \( MN = A \). Then \( N = M^{-1} \) and \( M \) is a finitely generated \( A \)-module. In particular, \( M \in I_A \).

**Proof.** By definition \( N \subset M^{-1} \), so \( N \subset M^{-1}MN \subset AN = N \). It follows that \( N = M^{-1} \). If \( M \) is finitely generated then \( M \in I_A \) by Remark 2 above. In any case, since \( A = MM^{-1} \), we have an equation \( 1 = \sum_{i=1}^{n} a_i b_i \) for some \( a_i \) in \( M \), some \( b_i \in M^{-1} \), and some \( n \). Then \( M = \sum_{i=1}^{n} Aa_i \). Indeed, if \( x \in M \), then \( x = \sum_i a_i(b_ix) \) lies in \( \sum_{i=1}^{n} Aa_i \). \( \square \)

We call a commutative ring a **semi-local ring** if it has only finitely many maximal ideals.

**Lemma 84.6.** If \( A \) is a semi-local domain, then \( Inv(A) = P_A \).

**Proof.** Let \( m_1, \ldots, m_n \) be the maximal ideals in \( A \) and suppose that \( M \in Inv(A) \), so \( MM^{-1} = A \). For each \( i = 1, \ldots, n \), choose \( a_i \in M \) and \( b_i \in M^{-1} \) such that \( a_i b_i \notin m_i \). Since \( \bigcap_{j \neq i} m_j \notin m_i \) for \( i = 1, \ldots, n \), there exist \( c_i \in (\bigcap_{j \neq i} m_j) \setminus m_i \) for \( i = 1, \ldots, n \). Let \( x = \sum c_i b_i \), an element of \( M^{-1} \). Therefore, \( xM \) is an ideal of \( A \). If \( xM = A \), then \( x \) is not zero and \( M = x^{-1}A \) is a principal fractional ideal. So suppose that \( xM < A \). Then \( xM \subset m_i \) for some \( i \). In particular, \( xa_i \in m_i \). As \( c_j b_j a_i \in m_i \) for every \( j \neq i \), we have \( c_i b_i a_i \) lies in \( m_i \). This is impossible. \( \square \)

Let \( A \) be a commutative ring and \( S \subset A \) a multiplicative set. If \( M \) is an \( A \)-module, then \( S^{-1}M := \{ \frac{m}{s} \mid s \in S, m \in M \} \) is an \( A \)- and an \( S^{-1}A \)-module in the obvious way [which is?]. If \( p \) is a prime ideal in \( A \), we denote by \( M_p \) the \( A_p \)-module \( S^{-1}M \) with \( S = A \setminus p \). The following is straight-forward:

**Lemma 84.7.** If \( A \) is a domain, \( S \) a multiplicative set in \( A \) not containing \( 0 \), and \( M \in Inv(A) \), then \( S^{-1}M \in Inv(S^{-1}A) \).

We have the following ‘local-global’ principle:

**Proposition 84.8.** Let \( A \) be a domain with quotient field \( F \) and \( M \) an \( A \)-submodule of \( F \). Then the following are equivalent:

1. \( M \in Inv(A) \).
2. \( M \) is a finitely generated \( A \)-module and \( M_p \in Inv(A) \), for all prime ideals \( p \) in \( A \).
(3) $M$ is a finitely generated $A$-module and $M_m \in \text{Inv}(A)$, for all maximal ideals $m$ in $A$.

(4) $M$ is a finitely generated $A$-module and $M_m \in P_A$, for all maximal ideals $m$ in $A$.

**Proof.** (1) $\Rightarrow$ (2) follows from lemmas 84.7 and 84.5, (2) $\Rightarrow$ (3) is immediate, and (3) $\Rightarrow$ (4) follows from Lemma 84.6.

(4) $\Rightarrow$ (1): Suppose that $MM^{-1} < A$. Then there exists a maximal ideal $m$ in $A$ with $MM^{-1} \subset m$. By assumption, $M$ is finitely generated, say $M = Aa_1 + \cdots + Aa_n$, and there exists a nonzero $x$ in $M$ satisfying $xA_m = M_m$. By the definition of localization, there exist nonzero $s_i$ in $A \setminus m$ satisfying $s_i a_i \in Ax$, $i = 1, \ldots, n$. Set nonzero $s = s_1 \cdots s_n$. Then $sM \subset Ax$, hence $sx^{-1}M \subset A$, i.e., $sx^{-1} \in M^{-1}$. It follows that $s = sx^{-1}x$ lies in $M^{-1}M \subset m$, a contradiction. Therefore, $MM^{-1} = A$ as desired. □

A domain is called a *Bézout domain* if every finitely generated ideal in $A$ is principal and a *Prüfer domain* if every finitely generated ideal is invertible.

**Remarks 84.9.**
1. Every Bézout domain is a Prüfer domain.
2. A domain $A$ is a Bézout domain if and only if $\text{Inv}(A) = P_A$.
3. A domain is a Prüfer domain if and only if every finitely generated fractional ideal is invertible.

We leave the following generalization of the fact that a DVR is a PID as an easy exercise:

**Proposition 84.10.** Let $A$ be a local domain. Then $A$ is a valuation ring if and only if $A$ is a Bézout domain.

The following shows that Prüfer domains generalize Dedekind domains when the Noetherian assumption is omitted.

**Theorem 84.11.** Let $A$ be a domain. Then the following are equivalent:
1. $A$ is a Prüfer domain.
2. $A_p$ is a valuation ring for all prime ideals $p$ in $A$.
3. $A_m$ is a valuation ring for all maximal ideals $m$ in $A$.

**Proof.** (1) $\Rightarrow$ (2): Let $\mathfrak{A}$ be a finitely generated ideal in $A_p$ with $p$ a prime ideal in $A$. Then there exists $a_i$ in $A \setminus p$ and $s_i$ in $S$, $i = 1, \ldots, n$, some $n$, such that $\mathfrak{A} = \sum_{i=1}^n A \frac{a_i}{s_i} = \left( \sum_{i=1}^n Aa_i \right)_p$. (Why?) Since $\sum Aa_i$ is invertible so is $\mathfrak{A}$. But then $\mathfrak{A}$ is principal by Proposition 84.8 as
$A_p$ is local. In particular, $A_p$ is Bézout, hence a valuation ring by the previous proposition.

(2) $\Rightarrow$ (3) is immediate.

(3) $\Rightarrow$ (1): Let $\mathfrak{A}$ be a finitely generated ideal in $A$. If $m$ is a maximal ideal in $A$, then $\mathfrak{A}_m$ is principal as $A_m$ is a valuation ring. Consequently, it is invertible by Proposition 84.8. □

We know that the integral closure of a Dedekind domain in a finite separable closure of its quotient field is also a Dedekind domain (although separable is not needed by the Krull-Akizuki Theorem 93.17 to be proven below). We shall generalize this to a Prüfer domain. Indeed we shall obtain the stronger result that the field extension need only be algebraic, i.e., we do not have to assume that it is finite nor separable. To do so we introduce a concept that leads to the ubiquity of valuation rings.

**Definition 84.12.** Let $A$ be a subring of $B$ and $\mathfrak{A} < A$ an ideal. We say that $\mathfrak{A}$ survives in $B$ if $\mathfrak{A}B < B$.

For example, by Claim 79.1, if $A$ is Dedekind domain then any ideal $\mathfrak{A} < A$ survives in $A_L$ for any field extension $L$ of $qf(A)$. Indeed this is true if $B/A$ is an integral extension for any commutative ring $A$ (as we shall prove in Theorem 90.14 below).

**Lemma 84.13. (Chevalley)** Let $A$ be a subring of the commutative ring $B$ and $\mathfrak{A} < A$ an ideal. If $u$ is a unit in $B$, then $\mathfrak{A}$ survives in either $A[u]$ or in $A[u^{-1}]$.

**Proof.** Suppose this is false. Then there exist elements $a_1, \ldots, a_n$ and $b_1, \ldots, b_m$ in $\mathfrak{A}$ satisfying:

(i) $1 = a_0 + a_1u + \cdots + a_nu^n$

(ii) $1 = b_0 + b_1u + \cdots + b_mu^{-m}$

in $B$, where we may assume that these have been chosen with $m, n \geq 0$ minimal. Since $\mathfrak{A} < A$, we may also assume that $m$ and $n$ are positive. We also assume that $m \leq n$. Note that $a_0 \neq 1$ and $b_0 \neq 1$ as $\mathfrak{A} < A$. Multiplying equation (ii) by $a_nu^n$ yields

(iii) $a_nu^n(1 - b_0) = b_1a_nu^{n-1} + \cdots + b_ma_nu^{n-m}$

and multiplying (i) by $1 - b_0$ yields

(iv) $1 - b_0 = (1 - b_0)a_0 + \cdots + (1 - b_0)a_nu^n$.

Plug (iii) into (iv) for $a_n(1 - b_0)u^n$ gives an equation of the form (i) of lesser degree in $u$ than $n$, a contradiction. □
**Theorem 84.14.** Let $A$ be a domain with quotient field $F$ with $K/F$ a field extension, and $\mathfrak{A} < A$ an ideal. Then there exists a valuation ring $B$ of $K$ containing $A$ with $\mathfrak{A}$ surviving in $B$.

**Proof.** Let

$$S = \{(A_\alpha, \mathfrak{A}_\alpha) \mid A \subset A_\alpha \subset K \text{ subrings, } \mathfrak{A} \subset \mathfrak{A}_\alpha < A_\alpha \text{ with } \mathfrak{A}_\alpha \text{ an ideal in } A_\alpha\}$$

Partially order $S$ by $\subseteq$ of pairs. Since $(A, \mathfrak{A}) \in S$, a Zorn Lemma argument produces a maximal element $(B, \mathfrak{B}) \in S$. Suppose that $B$ is not a valuation ring of $K$. Then there exists a nonzero element $x$ in $K$ with neither $x$ nor $x^{-1}$ in $B$. As $\mathfrak{B}$ survives in $B[x]$ or $B[x^{-1}]$, say $B[x]$, we would have $(B, \mathfrak{B}) < (B[x], B[x]\mathfrak{B})$. Since $B[x] \subset B[x, x^{-1}] \subset K$, the pair $(B[x], B[x]\mathfrak{B})$ lies in $S$, a contradiction.

**Definition 84.15.** Let $A$ and $B$ be local rings with maximal ideals $m, n$, respectively. We say that $B$ dominates $A$ if $A$ is a subring of $B$ and $m = A \cap n$.

**Corollary 84.16.** Let $A$ be a local domain with its quotient field $F$ lying in a field $K$. Set

$$D_K = \{B \text{ a local ring } \mid B \subset K\}$$

partially ordered by domination. Then

1. There exists a valuation ring $V$ in $D_K$ dominating $A$.
2. $A$ is a valuation ring in $D_F$ if and only if $A$ is maximal in $D_F$.

**Proof.** (1) follows from Theorem 84.14 with $\mathfrak{A}$ the maximal ideal of $A$.

(2): The proof of Theorem 84.14 implies the if statement. Conversely, suppose that $A \subset B \subset F$ are valuation rings of $F$ with maximal ideals $m_A$ and $m_B$, respectively, with $B$ dominating $A$, we must show that $A = B$. If $A < B$, then there exists a nonzero element $x$ in $B \setminus A \subset F$. Since $A$ is a valuation ring, $x^{-1}$ lies in $A \subset B$, so $x \in B^\times \setminus A$. It follows that $x^{-1}$ lies in $m_A = A \setminus A^\times$. Thus $m_B \cap A = m_A = A \setminus A^\times$, a contradiction. □

This corollary has an application to Prüfer domains.

**Corollary 84.17.** Let $R$ be a Prüfer domain with quotient field $F$. Suppose that $A$ is a valuation ring in $F$ containing $R$. Then there exists a prime ideal $p$ in $R$ such that $A = R_p$.

**Proof.** Let $m$ be the maximal ideal of $A$ and $p = m \cap R$, a prime ideal in $R$. Let $s \in R \setminus p$. Then $s^{-1}$ lies in $A$, lest $s \in m \cap R$. Therefore,
$R_p \subset A$. By Theorem 84.11 $R_p$ is a valuation ring, hence by Corollary 84.16, we have $R_p = A$.

This corollary gives a starting point for the theory of algebraic functions in one variable (geometrically curve theory). The algebraic formulation is the following:

**Corollary 84.18.** Let $F$ be a field and $A$ a valuation ring satisfying $F \subset A \subset F(t)$. Then there exists an irreducible polynomial $f$ in $F[t]$ satisfying $A = F[t]/(f)$ or $A = F[t^{-1}]/(t^{-1})$.

**Proof.** If $t$ lies in $A$, then $F[t] \subset A$. As $F[t]$ is a PID, it is Prüfer, so the result follows from Corollary 84.17. So we may assume that $t \notin A$. As $A$ is a valuation ring $t^{-1}$ lies in $A$. By Corollary 84.17, $A = F[t^{-1}]/p$ for some prime ideal in $F[t^{-1}]$. Since $t \notin A$, we have $t^{-1} \in p$, so $(t^{-1}) = p$.

We now wish to generalize Theorem 78.5 to the more general theorem about Prüfer domains previously mentioned. We need the following lemma.

**Lemma 84.19.** Let $A$ be an integrally closed local domain with quotient field $F$ and $a \in F$. Suppose that there exists a polynomial $f$ in $A[t]$ with $f(a) = 0$ and at least one coefficient of $f$ is a unit in $A$. Then either $a$ or $a^{-1}$ lies in $A$.

**Proof.** Suppose that $ba^n + ca^{n-1} + g(a) = f(a) = 0$ with $b, c \in A$ and $g \in A[t]$ satisfying $\deg g < n - 1$. If $b$ is a unit in $A$, then we are done as $A$ is integrally closed, so we may assume that $b \notin A^\times$. Then

(*) $0 = (ba)^n + c(ba)^{n-1} + b^{n-1}g(a) = (ba)^n + c(ba)^{n-1} + g_1(ba)$

for some $g_1 \in A[t]$ satisfying $\deg g_1 < n - 1$. Since $A$ is integrally closed, $ba$ lies in $A$, as $ba$ is integral over $A$ by (*). If $ba$ is a unit in $A$, then so is $b^{-1}a^{-1}$. This implies that $a^{-1} \in A$. Therefore, we may assume that $ba$ lies in the maximal ideal $m$ of $A$. We know that $(ba+c)a^{n-1} + g(a) = 0$. If $c \in A^\times$, then $ba + c \in A^\times$, as $A$ is local and $ba \in m$. But this means that $a$ is integral over $A$, so $a \in A$ as $A$ is integrally closed. Therefore, we may assume that $ba + c$ lies in $m$ and one of the coefficients of $g$ is a unit. As $\deg ((ba + c)t^{n-1} + g) < n$, we are done by induction on $n$.

**Theorem 84.20.** Let $A$ be a Prüfer domain with quotient field $F$. Let $L/F$ be an algebraic (possibly infinite) algebraic field extension. Then $A_L$ is Prüfer.
Proof. Let \( m \) be a maximal ideal in \( A_L \) and \( p = A \cap m \). We need to show that \( B_m \) is a valuation ring. Since \( A \) is a Pr"ufer domain, \( A_p \) is a valuation ring by Theorem 84.11. Let \( a \) be an element in \( L \) and \( f \) a nonzero polynomial in \( F[t] \) with \( f(a) = 0 \). Multiplying by a suitable element in \( A \), we may assume that \( f \in A[t] \). Since \( A_p \) is a valuation ring there exists a coefficient \( r \) of \( f \) dividing all the other coefficients of \( f \) in \( A \). Thus \( \frac{1}{r} f \in A_p[t] \) and has a coefficient that is a unit in \( A_p \), hence in \( (A_L)_m \). By the lemma, either \( a \) or \( a^{-1} \) lies in \( (A_L)_m \) as needed. \( \square \)

Corollary 84.21. Let \( A \) be a Dedekind domain with quotient field \( F \). Let \( L/F \) be an algebraic (possibly infinite) algebraic field extension. Then \( A_L \) is a Pr"ufer domain. In particular, if \( \bar{Q} \) is an algebraic closure of \( Q \), then \( \bar{Z}_{\bar{Q}} \) is a Pr"ufer domain.

Remarks 84.22. 1. If \( \bar{Q} \) is an algebraic closure of \( Z \), then \( \bar{Z}_{\bar{Q}} \) is not Noetherian so not a Dedekind domain.
2. The ring of entire functions on an open Riemann surface is a Pr"ufer domain. In fact, it is a B"ezout domain.
3. If in the proposition, B"ezout is substituted for Pr"ufer, then it can be shown that \( A_L \) is also B"ezout.
4. Bergman showed that if \( A = \mathbb{Z}[\sqrt{-5}] + t\mathbb{Q}[\sqrt{-5}][t] \), i.e., \( A \) is the ring of polynomials in \( \mathbb{Q}[\sqrt{-5}][t] \) having constant term in \( \mathbb{Z}[\sqrt{-5}] \), then \( A \) is not Noetherian, not B"ezout, but is Pr"ufer.
5. Every ideal in a Dedekind domain can be generated by two elements by Exercise 78.14(11). The same is false for finitely generated ideals in Pr"ufer domains. Schulting gave the first counterexample. It was subsequently generalized by Swan.

We give two further applications of valuation rings that are useful in commutative algebra. The first has as a consequence Zariski’s Lemma 40.10 that we saw was the key to the Hilbert Nullstellensatz in Section 40.

Theorem 84.23. Let \( F \) be an algebraically closed field and \( A \) a domain with \( R \subset A \) a subring. Suppose that \( A \) is a finitely generated \( R \)-algebra and \( a \in A \) nonzero. Then there exists an element \( 0 \neq r \in R \) satisfying the following condition: whenever \( \varphi : A \rightarrow F \) is a ring homomorphism with \( 0 \neq \varphi(r) \), there exists a ring homomorphism \( \psi : A \rightarrow F \) satisfying \( \psi|_R = \varphi \) and \( 0 \neq \psi(a) \).

Proof. By induction it suffices to assume that \( A = F[u] \).

Case 1. The element \( u \) is transcendental over \( R \).
Let \( a = \sum_{i=0}^{n} a_i u^i \) in \( F[u] \) with \( a_0 \) nonzero. We show \( r = a_0 \) works. So suppose that \( \phi : R \to F \) is a ring homomorphism with \( \phi(a_0) \neq 0 \). Since \( F \) is infinite, there exists an element \( x \in F \) such that \( \sum_{i=0}^{n} \phi(a_i) x^i \neq 0 \).

Let \( \psi : A \to F \) extend \( \phi \) by setting \( \psi(u) = x \). This works.

**Case 2.** The element \( u \) is algebraic over \( R \).

As \( R \) is a domain and \( a \) nonzero, \( a^{-1} \) is also algebraic over the quotient field of \( R \), so there exist equations:

\[
\begin{align*}
    a_0 u^n + a_1 u^{n-1} + \ldots + a_n &= 0 \\
    b_0 a^{-m} + b_1 a^{1-m} + \ldots + b_m &= 0
\end{align*}
\]

for appropriate \( a_i \) and \( b_j \) in \( A \) with \( a_0 \) and \( b_0 \) nonzero. It follows that \( u \) and \( a^{-1} \) are integral over \( R[(a_0 b_0)^{-1}] \). Set \( r = a_0 b_0 \) and suppose that \( \phi : R \to F \) is a ring homomorphism satisfying \( \phi(r) \neq 0 \). Set \( p = \ker \phi \).

Then \( p \) survives in \( R[r^{-1}] \) by choice, so there exists a valuation ring \( B \) in the quotient field of \( R[r^{-1}] \) such that \( pR[r^{-1}] \) survives in \( B \). If \( b \in B \), we can write \( b = r_1/r_2 \) with \( r_1, r_2 \in A \) and \( r_2 \notin p \) as \( qf(A) = qf(R[r^{-1}]) \). Hence we have a ring homomorphism \( \tilde{\phi} : B \to F \) given by \( r_1/r_2 \mapsto \phi(r_1)/\phi(r_2) \). We show that \( A \subset B \) and \( \psi : A \to F \) with \( \psi = \tilde{\phi}|_A \) work. Since the integral closure of \( R[r^{-1}] \) lies in \( B \), we have \( a^{-1} \) and \( u \) lie in \( B \). Thus \( A \subset B \) and \( a \in B^\times \). It follows that \( a \) does not lie in the maximal ideal of \( B \), so \( \tilde{\phi}(a) \neq 0 \).

**Corollary 84.24.** (Zariski’s Lemma) Let \( K/F \) be an extension of fields. Suppose that \( K \) is a finitely generated \( F \)-algebra. Then \( K/F \) is a finite field extension.

**Proof.** Let \( \bar{F} \) be an algebraic closure of \( F \). Then apply the theorem with \( R = F \), \( A = \bar{F} \), and \( a = 1 \). By the theorem, if \( x \in \bar{F} \) is transcendental over \( F \), the homomorphism \( \bar{F} \to \bar{F} \) sending \( x \to 0 \) cannot occur, which is false.

As a second application, we wish to show if \( A \) is an integrally closed domain, then so is \( A[t] \). To do this we first need the following:

**Theorem 84.25.** (Krull) Let \( A \) be a domain with quotient field \( F \) and

\[
S = \{ B \mid A \text{ is a subring of } B \text{ and } B \text{ is a valuation ring} \}.
\]

The \( A \) is integrally closed if and only if \( A = \bigcap_{S} B \).

**Proof.** (\( \Leftarrow \)): As valuation rings are integrally closed, this follows by Exercise 84.27(5).

(\( \Rightarrow \)): Certainly, \( A = \bigcap_{S} B \), so assume that there exists an element \( u \) in \( \bigcap_{S} B \setminus A \). Then \( u \) does not lie in \( A_f = A \), hence \( u \notin A[u^{-1}] \) by
Exercise 84.27(1). Let \( v = u^{-1} \), then \( vA[v] < A \), as \( v \notin A[v]^\times \). Hence there exists a valuation ring \( B_0 \subseteq F \) such that \( vA[v] \) survives in \( B_0 \). Since \( u \) lies in \( B \) for every \( u \in S \), we must have \( v \in B_0^\times \). But this means that \( vA[v] \) cannot survive in \( B_0 \), a contradiction \( \square \)

**Theorem 84.26.** Let \( A \) be an integrally closed domain. Then \( A[t] \) is integrally closed.

**Proof.** We have
\[
(*) \quad (\bigcap_{A \subseteq B \subseteq F} B)[t] = (\bigcap_{A \subseteq B \subseteq F} B[t]).
\]
(Why?) and by the previous theorem,
\[
A = \left( \bigcap_{A \subseteq B \subseteq F} B \right).
\]
As valuation rings are integrally closed by Exercise 84.27(5), we may assume that \( A \) is a valuation ring.

Let \( f \) be a nonzero element of \( F(t) \) (the quotient field of \( A[t] \)) be integral over \( A[t] \). We must show that \( f \) lies in \( A[t] \). Since \( F[t] \) is a UFD hence integrally closed, \( f \) lies in \( F[t] \). Write \( f = \frac{a}{g} \) with \( g \in A[t] \) and \( a \in A \) nonzero. If \( a \in A^\times \), we are done, so assume not. As \( A \) is a valuation ring, there exists a nonzero coefficient \( b \) of \( g \) dividing all the coefficients of \( g \), so we may write \( f = \frac{b}{a} h \) with \( h \in A[t] \) and \( h \) has some coefficient 1. If \( a | b \) we are done, so assume not. As \( A \) is a valuation ring, we must have \( b | a \) in \( A \), so \( f = \frac{h}{a} \) for some nonzero element \( c \) in \( A \). As \( c \) is not a unit in \( A \) and \( A \) is local it lies in the maximal ideal \( m \) of \( A \). As \( f \) is integral over \( A[t] \), there exists polynomials, \( \alpha_1(t), \ldots, \alpha_n(t) \) in \( A[t] \) for some \( n \) satisfying \( f^n + \alpha_1(t)f^{n-1} + \cdots + \alpha_n(t) = 0 \) in \( F[t] \). Hence \( h^n + c\alpha_1(t)h^{n-1} + \cdots + c^n\alpha_n(t) = 0 \) in \( A[t] \). Consequently, \( h^n \equiv 0 \) mod \( A[t]m \). But \( A[t]/mA[t] \) is a domain, \( \text{i.e., } A[t]/m \) is a prime ideal in \( A[t] \), so \( h \in A[t]/mA[t] \) (why?). This implies that every coefficient of \( h \) lies in \( m \), a contradiction. \( \square \)

**Exercises 84.27.**

1. Let \( A \) be a commutative ring and \( u \) a unit in \( A \). If \( R \) is a subring of \( A \), show that \( u^{-1} \) is integral over \( R \) if and only if \( u^{-1} \) lies in \( R[u] \).
2. Let \( A \) be a commutative ring. Then \( A \) is a domain if and only if \( A_m \) is a domain for all maximal ideals \( m \) of \( A \).
3. Show 84.7.
4. Establish Proposition 84.10
5. Prove that the intersection of integrally closed domains each of which lies in a common field is integrally closed.
In the previous chapter we studied the algebra of Dedekind domains. Since the ring of integral elements (i.e., the ring of algebraic integers) in an (algebraic) number field (a finite field extensions of the rational numbers) is a Dedekind domain, we also developed some of the properties of a ring of algebraic integers in that chapter. For example, we investigated Hilbert ramification theory of a ring of algebraic numbers including a discussion of its discriminant. We also fully computed the theory for quadratic number rings. In the case of a quadratic number ring, the units were easy to compute (especially in the non-real case). In this chapter, we wish to classify the units in any ring of algebraic numbers. We also wish to investigate the discriminant of a ring of algebraic numbers more thoroughly. Two of the main advantages in the special case of number theory is that its ideals are finitely generated free $\mathbb{Z}$-modules and its quotient fields by its nonzero prime ideals are finite fields.

85. Ideal and Counting Norms

In this section, we need a preliminary discussion of norms, that we could have done previously. We know that if $A$ is a integrally closed domain with quotient field $F$ and $K$ is a finite separable extension, that $N_{K/F} : K \to F$ maps $A_K$ to $A$. If we also assume that $A$ is a Dedekind domain, then we have unique factorization of fractional ideals into products of prime ideals in the Dedekind domain $A_K$. This will allow us to define the norm of a fractional ideal of $A_K$ to be a fractional ideal in $A$. If $F = \mathbb{Z}$, i.e., $A_K = \mathbb{Z}_K$, then the norm of an element in $\mathbb{Z}_K$ is an integer. Moreover, if $p$ is a nonzero prime ideal in $\mathbb{Z}_K$, then $\mathbb{Z}_K/p$ is a finite (separable) extension of the finite field $\mathbb{Z}/p \cap \mathbb{Z}$ of degree $f(p/p \cap \mathbb{Z})$. This allows us to define another norm. We will then show how these norms are interrelated.

Since we have investigated the case for an arbitrary Dedekind domain, we shall begin by studying the case of norms of ideals in this more general setting.

We shall need the following that we leave as an exercise.
Lemma 85.1. Let $K/F$ be a finite separable field extension and $A$ a Dedekind domain with quotient field $F$. If $\mathfrak{A} \in I_A$, then $\mathfrak{A}_K \cap A = \mathfrak{A}$ and $I_A \rightarrow I_{A_K}$ by $\mathfrak{A} \rightarrow \mathfrak{A} A_K$ is a group monomorphism.

Definition 85.2. Let $K/F$ be a finite separable field extension and $A$ a Dedekind domain with quotient field $F$. If $\mathfrak{P}$ is a nonzero prime ideal in $A_K$, we know that $f(\mathfrak{P}/\mathfrak{P} \cap A) := \dim_{A/(\mathfrak{P} \cap A)}(A_K/\mathfrak{P})$ is a finite integer by Theorem 79.2. By the unique factorization of fractional ideals in a Dedekind domain, we can define $N_{K/F} : I_{A_K} \rightarrow I_A$ to be the group homomorphism induced by $\mathfrak{P} \mapsto (\mathfrak{P} \cap A)^{f(\mathfrak{P}/\mathfrak{P} \cap A)}$ called the ideal norm map.

By the multiplicativity the ideal norm, we have:

Remark 85.3. Let $K/F$ be a finite separable field extension and $A$ a Dedekind domain with quotient field $F$. If $\mathfrak{A} \in I_{A_K}$, then $N_{K/F}(\mathfrak{A} A_K) = \mathfrak{A}^{[K:F]}$, i.e., the ideal generated by norms of elements.

We now show that ideal norms have similar properties as norms of elements and how the two are interrelated. We need two lemmas.

Lemma 85.4. Let $K/F$ be a finite separable field extension and $A$ a Dedekind domain with quotient field $F$. If $\mathfrak{A} \in I_A$, then $N_{K/F}(\mathfrak{A} A_K) = \mathfrak{A}^{[K:F]}$.

Proof. As the ideal norm is a group homomorphism, we may assume that $\mathfrak{A} = \mathfrak{p}$ is a nonzero prime ideal in $A$. Let $\mathfrak{p} A_K = \mathfrak{P}_1^{e_1} \cdots \mathfrak{P}_r^{e_r}$ be a prime factorization of $\mathfrak{p} A_K$ in $A_K$ and $f_i = f(\mathfrak{P}_i/\mathfrak{p})$ for $i = 1, \ldots, r$. By definition, $N_{K/F}(\mathfrak{p} A_K) = \mathfrak{p}^{e_1 f_1} \cdots \mathfrak{p}^{e_r f_r}$, hence $N_{K/F}(\mathfrak{p} A_K) = \mathfrak{p}^{[K:F]}$ by Theorem 79.2.

Lemma 85.5. Let $L/F$ be a finite Galois extension and $A$ a Dedekind domain with quotient field $F$, $\mathfrak{p} A_L = \mathfrak{P}_1^{e_1} \cdots \mathfrak{P}_r^{e_r}$ a factorization, with $\mathfrak{p}$ a nonzero prime in $A$. Then $e = e_i$ and $f_i = f(\mathfrak{P}_i/\mathfrak{p})$ for $i = 1, \ldots, r$ and $N_{L/F}(\mathfrak{P}_i) A_L = (\mathfrak{P}_1 \cdots \mathfrak{P}_r)^{e_i f} = \prod_{G(L/F)} \sigma(\mathfrak{P}_i)$.

Proof. By Proposition 80.2 (and Proposition 80.1), we have $e = e_i$ and $f = f_i$ for $i = 1, \ldots, r$ and $G(L/F)$ acts transitively on $\mathfrak{P}_1, \ldots, \mathfrak{P}_r$. The result follows by Lemma 85.4.

Proposition 85.6. Let $K/F$ be a finite separable field extension and $A$ a Dedekind domain with quotient field $F$ and $x \in K^\times$. Then $N_{K/F}(x) A = N_{K/F}(x A_K)$, where the norm on the right-hand side is the ideal norm and the norm on the left-hand side is the field norm.
85. IDEAL AND COUNTING NORMS

Let $L/F$ be the normal closure of $K/F$ and $x \in K^\times$, then

$$N_{L/F}(x)A = N_{K/F}(x^{[L:K]})A = (N_{K/F}(x)A)^{[L:K]}$$

$$N_{L/F}(xA_L) = N_{K/F}((xA_K)^{[L:K]}) = (N_{K/F}(xA_K)^{[L:K]})$$

by Lemma 85.4 and the multiplicativity of $f$ on towers by Exercise 79.7(2). By uniqueness of factorization of ideals, we are reduced to the case that $L = K$, i.e., $K/F$ is Galois. We then have

$$N_{K/F}(x)A_K = \prod_{G(K/F)} \sigma(x)A_K = \prod_{G(K/F)} \sigma(xA_K)$$

$$= \prod_{G(K/F)} \prod_{\mathfrak{p}} \sigma(\mathfrak{p}^{v_{\mathfrak{p}}(x)}) = \prod_{G(K/F)} \prod_{\mathfrak{p}} (\mathfrak{p})^{v_{\mathfrak{p}}(x)}$$

$$= N(xA_K)A_K$$

by Lemma 85.5 and the multiplicativity of $N$. By Lemma 85.1, $\mathfrak{A}A_K \cap A = \mathfrak{A}$, so the result follows.

We also need another norm in the case of a number field.

**Definition 85.7.** Let $K$ be a number field. If $\mathfrak{A}$ is a nonzero ideal in $\mathbb{Z}_K$, define the **counting norm** of $\mathfrak{A}$ to be the integer $N(\mathfrak{A}) := |\mathbb{Z}_K/\mathfrak{A}|$.

We want to extend the counting norm $\mathfrak{N}$ to a group homomorphism $\mathfrak{N} : I_{\mathbb{Z}_K} \rightarrow \mathbb{Q}^\times$. To do so, we now show that the counting norm is multiplicative. We also show how it is related to the ideal norm.

**Proposition 85.8.** Let $K$ be a number field and $\mathfrak{A}, \mathfrak{B} \subset \mathbb{Z}_K$ nonzero ideals. Then

1. $\mathfrak{N}(\mathfrak{A} \cup \mathfrak{B}) = \mathfrak{N}(\mathfrak{A}) \mathfrak{N}(\mathfrak{B})$.
2. $\mathfrak{N}(\mathfrak{A} \mathfrak{B}) = \mathfrak{N}(\mathfrak{A}) \mathfrak{N}(\mathfrak{B})$.
3. If $E/K$ is a finite extension, then $\mathfrak{N}(\mathfrak{A} \mathfrak{Z}_E) = \mathfrak{N}(\mathfrak{A} \mathfrak{Z}_K)^{[K:F]}$.
4. If $a \in \mathbb{Z}_K$ is nonzero, then $\mathfrak{N}_{K/Q}(a \mathfrak{Z}_K) = |N_{K/Q}(a)|$.
5. $\mathfrak{N}(\mathfrak{A}) = |N_{K/Q}(\mathfrak{A})|$, where we mean a positive generator for the ideal $N_{K/Q}(\mathfrak{A})$ on the right-hand side.

**Proof.**

1. If $\mathfrak{A}$ and $\mathfrak{B}$ are relatively prime, we have $\mathbb{Z}_K = \mathfrak{A} + \mathfrak{B}$. In this case, we have $\mathfrak{N}(\mathfrak{A} \mathfrak{B}) = \mathfrak{N}(\mathfrak{A}) \mathfrak{N}(\mathfrak{B})$ by the Chinese Remainder Theorem. So we need only show:

   **Claim** Let $\mathfrak{p}$ be a nonzero prime ideal in $\mathbb{Z}_K$. Then $\mathfrak{N}(\mathfrak{p}^r) = \mathfrak{N}(\mathfrak{p})^r$ for all positive integers $r$:

   Fix $r$ and choose $\alpha \in \mathfrak{p}^r \setminus \mathfrak{p}^{r+1}$. Then we have (cf. Exercise 79.14(1))

   1. $\mathfrak{p}^r = \alpha \mathbb{Z}_K + \mathfrak{p}^{r+1}$, the greatest common divisor of $\alpha \mathbb{Z}_K$ and $\mathfrak{p}^{r+1}$.
(ii) \( \alpha p^r = \alpha \mathbb{Z}_K \cap p^{r+1} \), the least common multiple of \( \alpha \mathbb{Z}_K \) and \( p^{r+1} \).

Since \( \alpha p^r \subset p^{r+1} \), it follows that for each \( r \), we have a \( \mathbb{Z}_K \)-isomorphism \( \varphi : \mathbb{Z}_K / p \to p^r / p^{r+1} \) given by \( \beta \mapsto \beta \alpha \), which is also an isomorphism of one dimensional \( \mathbb{Z}_K / p \)-vector spaces. If follows (say by the Jordan-Hölder Theorem) that

\[
\mathfrak{N}(p^s) = |\mathbb{Z}_K / p^s| = \prod_{i=0}^{s-1} |p^i / p^{i+1}| = \mathfrak{N}(p)^s
\]

where \( p^0 = \mathbb{Z}_K \). This proves the Claim and establishes (2).

(1) follows from (2).

(3): By (2), we may assume that \( \mathfrak{A} = p \) is a prime ideal in \( \mathbb{Z}_K \). Let \( p \mathbb{Z}_E = \mathfrak{P}_1^{e_1} \cdots \mathfrak{P}_r^{e_r} \) be a factorization in \( \mathbb{Z}_E \) and \( f_i = f(\mathfrak{P}_i / p) \) for \( i = 1, \ldots, r \). By (2), we have \( \mathfrak{N}(\mathfrak{P}_i^{e_i}) = \mathfrak{N}(\mathfrak{P}_i)^{e_i} = \mathfrak{N}(p)^{e_if_i} \), hence

\[
\mathfrak{N}(p \mathbb{Z}_E) = \mathfrak{N}(p)^{\sum_{i=1}^r e_if_i} = \mathfrak{N}(p)^{[K:F]}
\]

by Theorem 79.2 as needed.

(5): As above, we may assume that \( \mathfrak{A} = p \) is a prime ideal in \( \mathbb{Z}_K \). Let \( p \mathbb{Z} = p \cap \mathbb{Z} \). Then \( \mathfrak{N}(p) = p^f \) with \( f = [\mathbb{Z}_K / p : \mathbb{Z} / p \mathbb{Z}] \). By definition \( N_{\mathbb{K}/\mathbb{Q}}(p) = p^f \mathbb{Z} \). Statement (5) now follows.

(4) follows from (5).

□

**Corollary 85.9.** Let \( K \) be a number field and \( \mathfrak{A}, \mathfrak{B} \subset \mathbb{Z}_K \) nonzero ideals. Then \( \mathfrak{N} \) extends to a group homomorphism \( \mathfrak{N} : I_{\mathbb{Z}_K} \to \mathbb{Q}^* \).

Moreover, the counting norm satisfies all the properties in Proposition 85.8 for all \( \mathfrak{A}, \mathfrak{B} \in I_A \).

The key property that we shall need, besides the multiplicativity of the counting norm, is that we can replace the ideal norm of a principal fractional ideal by the absolute value of the field norm of a positive generator of it. In particular, if \( u \) is a unit in \( \mathbb{A}_K \) and \( \alpha \in \mathbb{Z}_K \), then

\[
\mathfrak{N}(a \mathbb{Z}_K) = |N_{\mathbb{K}/\mathbb{Q}}(au)|.
\]

**Exercises 85.10.**

Let \( K/F \) be a finite separable field extension and \( A \) a Dedekind domain with quotient field \( F \). Suppose that \( \mathfrak{A} \) and \( \mathfrak{B} \) are ideals in \( A \). Prove the following:

1. If \( \mathfrak{A} | \mathfrak{A}_K \) in \( A_K \), then \( \mathfrak{A} | \mathfrak{B} \) in \( A \).
2. Lemma 85.1.
3. Remark 85.3
86. LATTICES IN NUMBER FIELDS

In this section, we show how to view a ring of algebraic integers as a lattice in an appropriate Euclidean space, i.e., a finite dimensional real inner product space in an appropriate way and thereby establish that it is a discrete subgroup there. Our goal in this section is to prove the Minkowski Lattice Point Theorem, a most useful result in number theory. We shall assume the necessary topology and analysis in this section.

We start with the setup that we shall need throughout the rest of this chapter. Let $K$ be a number field of degree $n$ over $\mathbb{Q}$. Then there exist $n$ embeddings of $K$ into $\mathbb{C}$ by Proposition 55.11 and the comment following (cf. Proposition 57.2).

We write these embeddings as $\sigma_1, \ldots, \sigma_r : K \to \mathbb{C}$ for all the real embeddings (i.e., those with image in $\mathbb{R}$) and $\tau_1, \ldots, \tau_r : K \to \mathbb{C}$ for the complex embeddings (i.e., those with image not contained in $\mathbb{R}$), where $\bar{\sigma}$ is complex conjugation and $\sigma_{r_1+j}$ is the composition $\bar{\sigma} \circ \sigma_{r_1+j}$ for $i = 1, \ldots, r_2$. In particular, we have $[K : \mathbb{Q}] = n = r_1 + 2r_2$.

If $x \in K$, we shall use the following notation:

$$x^{(i)} = \sigma_i(x), \quad i = 1, \ldots, n,$$

where $\sigma_{r_1+j} = \sigma_{r_1+j}$ for $j = 1, \ldots, r_2$.

We view $\mathbb{R}^{r_1} \times \mathbb{C}^{r_2}$ as a real vector space (with $r_1 = 0$ or $r_2 = 0$ allowed and with each $\mathbb{C}$ having basis $\{1, \sqrt{-1}\}$).

Therefore, if $V = \mathbb{R}^{r_1} \times \mathbb{C}^{r_2}$, we have an $\mathbb{Z}$-embedding

$$i_{\mathbb{C}} : K \to \mathbb{R}^{r_1} \times \mathbb{C}^{r_2} \text{ via } x \mapsto (\sigma_1(x), \ldots, \sigma_{r_1+r_2}(x)) = (x^{(1)}, \ldots, x^{(r_1+r_2)}),$$

where the map will usually be suppressed. If $V = \mathbb{R}^n$, $n = r_1 + 2r_2$, we shall also need another additive embedding given by

$$i_{\mathbb{R}} : K \to \mathbb{R}^n \text{ where } x \mapsto (x^{(1)}, \ldots, Re(x^{(r_1+1)}), Im(x^{(r_1+1)}), \ldots, Re(x^{(r_1+r_2)}), Im(x^{(r_1+r_2)})).$$

**Definition 86.1.** Suppose that $V$ is an $n$-dimensional real inner product space, e.g., $\mathbb{R}^{r_1} \times \mathbb{C}^{r_2}$ above, and $\{v_1, \ldots, v_r\}$ a linearly independent set in $V$. Let $\mathcal{L} := \bigoplus_{i=1}^r \mathbb{Z}v_i$. Then $\mathcal{L}$ is a free $\mathbb{Z}$-module of rank $r$ called an $r$-dimensional lattice in $V$ with lattice basis $\{v_1, \ldots, v_r\}$. We
often write $\mathcal{L}$ as $\mathcal{L}(v_1, \ldots, v_r)$ for clarity, when we wish to highlight the lattice basis. If $r = n = \dim V$, we call $\mathcal{L}$ a full lattice in $V$, i.e., if \{v_1, \ldots, v_r\} is a basis for $V$. If $\mathcal{L}$ is a full lattice in $V$, we let
\[
D_\mathcal{L} = D_{\mathcal{L}(v_1, \ldots, v_n)} := \left\{ \sum_{i=1}^{n} a_i v_i \mid 0 \leq a_i < 1, a_i \in \mathbb{R}, 1 \leq i \leq n \right\}
\]
called the fundamental domain (or parallelepiped) of $\mathcal{L}$.

Note that the fundamental domain $D_\mathcal{L}$ of $\mathcal{L}$ depends on \{v_1, \ldots, v_n\} even though $\mathcal{L}$ does not.

**Remarks 86.2.** Let $V$ be a finite dimensional real inner product space.
1. The determinant of the change of basis matrix of any two lattices bases of a lattice $\mathcal{L}$ in $V$ is $\pm 1$.
2. If $\mathcal{L}$ is a full lattice in $V$, then $V/\mathcal{L}$ is a compact space.
3. Suppose that $\mathcal{L} = \mathcal{L}(v_1, \ldots, v_n)$ is a full lattice in $V$ and $\mathcal{B} = \{e_1, \ldots, e_n\}$ is an orthonormal basis for $V$. If $v_i = \sum_{i=1}^{n} a_{ij} e_j$, then
\[
\text{Vol}(D_\mathcal{L}) = \text{Vol}(V/\mathcal{L}) = |\det(a_{ij})|
\]
and is independent of lattice basis by (1).

**Example 86.3.** Let $V = \mathbb{R}^n \times \mathbb{C}^r$ or $\mathbb{R}^n$, $K$ a number field of degree $n$. We know that $\mathbb{Z}_K$ is a free abelian group of rank $n$, so is a full lattice when viewed in the image of $K$. Let $0 < \mathfrak{a} < \mathbb{Z}_K$ be an ideal, hence also $\mathbb{Z}$-free. By Corollary 78.6, there exists a nonzero $m \in \mathfrak{a} \cap \mathbb{Z}$, so $0 < m\mathbb{Z}_K \subset \mathfrak{a} < \mathbb{Z}_K$. In particular, $\mathbb{Z}_K/m\mathbb{Z}_K$ is a finite additive torsion group. Since $\mathbb{Z}_K/\mathfrak{a}$ is a finite group, $\mathfrak{a}$ is also a full lattice in $V$. More generally, if $\mathfrak{a}$ is a fractional ideal, there exists a nonzero $c \in \mathbb{Z}_K$ such that $c\mathfrak{a}$ is an ideal in $\mathbb{Z}_K$. Therefore, $N_{K/\mathbb{Q}}(c)$ is an integer so $N_{K/\mathbb{Q}}(c)\mathfrak{a} \subset \mathbb{Z}_K$ also. As $N_{K/\mathbb{Q}}(c) \in \mathbb{Z}$, we also have $[N_{K/\mathbb{Q}}(c)^{-1}\mathbb{Z}_K : \mathfrak{a}]$ is finite and $\mathfrak{a}$ is a full lattice in $V$. Note under $i_K$ and $i_{\mathfrak{a}K}$, these induce lattices in $\mathbb{R}^{r_1} \times \mathbb{C}^{r_2}$ and $\mathbb{R}^{r_1+2r_2}$, respectively, with fundamental domains having the same volume. In particular, this will also be true for the lattices induced by any $\mathfrak{a} \in I_{\mathbb{Z}_K}$

**Lemma 86.4.** Let $V$ be a finite dimensional real inner product space and $\mathcal{L} = \mathcal{L}(v_1, \ldots, v_n)$ be a full lattice in $V$. Then $\{\lambda + D_\mathcal{L} \mid \lambda \in \mathcal{L}\}$ partitions $V$.

**Proof.** Let $v \in V$, then $V = \sum_{i=1}^{n} (q_i + r_i)v_i$ with $0 \leq r_i < 1$ and $q_i \in \mathbb{Z}$ for $i = 1, \ldots, n$. Thus the sets $\lambda + D_\mathcal{L}$ cover $V$. If we have $\langle \lambda + D_\mathcal{L} \rangle \cap (\lambda' + D_\mathcal{L}) \neq \emptyset$, then $D_\mathcal{L} \cap ((\lambda - \lambda') + D_\mathcal{L}) \neq \emptyset$. It follows that $\lambda = \lambda'$.

$\square$
If \( V \) is a finite dimensional real inner product space and \( \mathcal{L} \) an additive subgroup of \( V \), we call \( \mathcal{L} \) a **discrete subgroup** of \( V \) if \( \mathcal{L} \) meets every compact subset of \( V \) in a finite set.

**Proposition 86.5.** Let \( V \) be a finite dimensional real inner product space and \( \mathcal{L} \) an additive subgroup of \( V \). Then \( \mathcal{L} \) is a lattice if and only if \( \mathcal{L} \) is discrete. In particular, if \( \mathcal{L} \) is discrete, then it is a finitely generated additive group.

**Proof.** (\( \Rightarrow \)): Let \( C \subset V \) be a compact set, \( \mathcal{L} = \mathcal{L}(w_1, \ldots, w_r) \).

If \( \sum_{i=1}^r a_i w_i \) lies in \( C \cap \mathcal{L} \), then the integers \( |a_i| \)'s must be uniformly bounded.

(\( \Leftarrow \)): Let \( \{v_1, \ldots, v_n\} \) be a basis for \( V \). We may assume that \( \mathcal{L} \neq 0 \), and we induct on \( n = \dim V \). If \( n = 1 \), choose \( a > 0 \) minimal such that \( v = av_1 \) lies in \( \mathcal{L} \). (Such an \( a \) exists as we can intersect \( \mathcal{L} \) with a large enough ball centered at the origin.) Let \( bv_1 \in \mathcal{L} \) with \( b \neq 0 \). Write \( b = da + c \) with \( 0 \leq c < a \) and \( d \in \mathbb{Z} \). By hypothesis, we must have \( c = 0 \), so \( bv_1 = dav_1 \), hence \( \mathcal{L} = \mathbb{Z} v \).

Now suppose that \( n > 1 \). We may assume that \( \mathcal{L} \) lies in no proper subspace of \( V \) by induction. Set \( V_0 = \bigcap_{i=1}^{n-1} \mathbb{R} v_i \) and \( \mathcal{L}_0 = \mathcal{L} \cap V_0 < \mathcal{L} \). So \( \mathcal{L}_0 \) cannot lie in any proper subspace of \( V_0 \), hence by induction \( \mathcal{L}_0 \) is a full lattice of rank \( n-1 \) in \( V_0 \). Let \( \mathcal{L}_0 = \mathcal{L}(u_1, \ldots, u_{n-1}) \).

Then \( \{u_1, \ldots, u_{n-1}, v_n\} \) is a basis for \( V \). Since \( \mathcal{L}_0 < \mathcal{L} \), we see that there exists an element \( \lambda \in \mathcal{L} \) satisfying \( \lambda = \sum_{i=1}^{n-1} r_i u_i + r_n v_n \) with \( 0 \leq r_i < 1, i = 1, \ldots, n-1 \), and \( r_n > 0 \). Let

\[
C = \left\{ \sum_{i=1}^{n-1} x_i u_i + x_n v_n \mid 0 \leq x_i \leq 1, i = 1, \ldots, n-1, |x_n| \leq r_n \right\},
\]

a compact subset of \( V \). Hence \( \mathcal{L} \cap C \) is a finite set. Choose \( u_n \in \mathcal{L} \setminus \mathcal{L}_0 \) to satisfy

\[
u_n = \sum_{i=1}^{n-1} x_i u_i + x_n v_n, \quad 0 \leq x_i < 1, i = 1, \ldots, n-1, \quad |x_n| > 0 \text{ minimal.}
\]

**Claim.** \( \mathcal{L} = \mathcal{L}(u_1, \ldots, u_n) \).

If not then there exists a \( w \in \mathcal{L} \setminus \mathcal{L}(u_1, \ldots, u_n) \). Since \( \mathcal{L}(u_1, \ldots, u_n) = \mathcal{L}(u_1, \ldots, u_{n-1}) \oplus \mathbb{Z} u_n \), it is a full lattice in \( V \). Subtracting a suitable nonzero integral multiple of \( u_n \) and then by a suitable element of \( \mathcal{L}_0 \), we see, using \( \mathbb{Z} u_n \cap \mathcal{L}_0 = 0 \), that there exist \( 0 \leq b_i < 1, i = 1, \ldots, n-1 \), and \( 0 < a_n |x_n| < |x_n| \) satisfying \( w_0 = \sum_{i=1}^{n-1} b_i u_i + a_n x_n v_n \) lies in \( \mathcal{L} \).

This contradicts the minimality of \( |x_n| \). (Cf. the \( n = 1 \) case.) \( \square \)
Lemma 86.6. Let $V$ be a finite dimensional real inner product space and $L_1 \subset L_2$ be full lattices in $V$. Then $L_2/L_1$ is finite and
\[ \mu(D_{L_1}) = \text{Vol}(V/L_1) = [L_2 : L_1] \text{Vol}(V/L_2) = [L_2 : L_1] \mu(D_{L_2}). \]

Proof. Since $L_2/L_1$ is a finitely generated group of rank zero, it is torsion hence finite. Let $L_2 = L(v_1, \ldots, v_n)$. By linear algebra, there exists a $\mathbb{Z}$-basis $\{w_1, \ldots, w_n\}$ for $L_1$ satisfying
\[ (i) \quad w_i = c_{ii} v_i + c_{i,i+1} v_{i+1} + \cdots + c_{in} v_n, \quad c_{ij} \in \mathbb{Z} \text{ and } c_{ii} > 0 \text{ for all } i,j. \]
Moreover,
\[ \alpha := |\det(c_{ij})| = c_{11} \cdots c_{nn} \text{ is independent of } \{w_1, \ldots, w_n\} \]
(as change of bases matrices for lattices lie in $\text{GL}_n(\mathbb{Z})$). So $\text{Vol}(V/L_1) = \alpha \text{Vol}(V/L_2)$. We finish by proving:

Claim. $\alpha = [L_2 : L_1]$;

Let $v \in L_2$. Using $(i)$, we see that there exist $x_i \in \mathbb{Z}$, $0 \leq x_i < c_{ii}$, satisfying
\[ (ii) \quad v \equiv x_1 v_1 + \cdots + x_n v_n \mod L_1. \]
It follows that $\alpha \geq [L_2 : L_1]$. (as the number of possible $x_i$ is $c_{ii}$ all $i$.)

To show that $\alpha = [L_2 : L_1]$, it suffices to show that all the $v$’s in $(ii)$ give different cosets, i.e., given
\[ v = \sum_{i=1}^{n} x_i v_i, \quad v' = \sum_{i=1}^{n} x'_i v_i \text{ in } L_2 \text{ and } v \equiv v' \mod L_1 \]
with $x_i, x'_i \in \mathbb{Z}$ satisfying $0 \leq x_i, x'_i < c_{ii}$,
we must show $v = v'$. But if this is the case, $v - v' = \sum_{j=1}^{n} b_j w_j$ for some $b_j \in \mathbb{Z}$, hence
\[ (iii) \quad v - v' = \sum_{i=1}^{n} (x_i - x'_i) v_i = \sum_{j=1}^{n} b_j w_j = \sum_{j=1}^{n} b_j (\sum_{i=1}^{n} c_{ji} v_i) \]
by $(i)$. If $v \neq v'$, then there exists a minimal $s$ satisfying $x_s \neq x'_s$. Then by $(iii)$, we have $b_j = 0$ for all $j < s$ and $c_{ss} b_s = x_s - x'_s \neq 0$. However, $0 < |x_s - x'_s| < c_{ss} \leq |c_{ss} b_s|$, a contradiction. So $v = v'$ and the claim, hence the lemma, is proven. \qed

The lemma allows us to establish the following important computation.
Proposition 86.7. Let $K$ be a number field of degree $n$, $V = \mathbb{R}^n$ (or $\mathbb{R}^{r_1} \times \mathbb{C}^{r_2}$), and $\mathfrak{A} \in I_{Z_K}$. Then

$$\mu(D_\mathfrak{A}) = \text{Vol}(V/\mathfrak{A}) = \mathfrak{N}(\mathfrak{A})2^{-r_2}\sqrt{|d_K|}.$$ 

Proof. Let $\{w_1, \ldots, w_n\}$ be a basis for $K$ and $V = \mathbb{R}^{r_1} \times \mathbb{C}^{r_2}$. Write $w_i^{(r_1+j)} = x_i^{(r_1+j)} + \sqrt{-1}y_i^{(r_1+j)}$, $x_i^{(r_1+j)}$, $y_i^{(r_1+j)} \in \mathbb{R}$ for $j = 1, \ldots, r_2$. Set $\Delta = \Delta_{K/\mathbb{Q}}(w_1, \ldots, w_n)$ and $\mathcal{L} = \mathcal{L}(w_1, \ldots, w_n)$. Then we have

$$\Delta = \text{det} \begin{pmatrix}
  w_1^{(1)} & \cdots & w_1^{(1)} \\
  \vdots & & \vdots \\
  x_1^{(r_1+1)} + \sqrt{-1}y_1^{(r_1+1)} & \cdots & x_n^{(r_1+1)} + \sqrt{-1}y_n^{(r_1+1)} \\
  \vdots & & \vdots \\
  x_1^{(r_1+1)} - \sqrt{-1}y_1^{(r_1+1)} & \cdots & x_n^{(r_1+1)} - \sqrt{-1}y_n^{(r_1+1)}
\end{pmatrix}^2$$

$$= (-2\sqrt{-1})^{2r_2} \text{det} \begin{pmatrix}
  w_1^{(1)} & \cdots & w_n^{(1)} \\
  \vdots & & \vdots \\
  x_1^{(r_1+1)} & \cdots & x_n^{(r_1+1)} \\
  \vdots & & \vdots \\
  y_1^{(r_1+1)} & \cdots & y_n^{(r_1+1)}
\end{pmatrix}^2$$

$$= (-1)^{r_2}2^{2r_2} \text{Vol}(V/\mathcal{L})^2.$$ 

Therefore, (computing $\text{Vol}(V/\mathcal{L})$ using the standard basis), we have

$$\text{Vol}(V/\mathcal{L}) = 2^{-r_2}\sqrt{|\Delta|}.$$ 

If $\{w_1, \ldots, w_n\}$ is an integral basis for $Z_K$, then $\text{Vol}(V/Z_K) = 2^{-r_2}\sqrt{|d_K|}$. If $0 < \mathfrak{A} \subset Z_K$ is an ideal, by Lemma 86.6,

$$\text{Vol}(V/\mathfrak{A}) = [Z_K : \mathfrak{A}]2^{r_2}\sqrt{|d_K|} = \mathfrak{N}(\mathfrak{A})2^{r_2}\sqrt{|d_K|}.$$ 

Finally, suppose that $\mathfrak{A} \in I_{Z_K}$. By Remark 76.14(1), there exists a nonzero integer $m$ satisfying $m\mathfrak{A} \subset \mathfrak{A}$. Then

$$m^n\text{Vol}(V/\mathfrak{A}) = \text{Vol}(V/m\mathfrak{A}) = \mathfrak{N}(m\mathfrak{A})2^{r_2}\sqrt{|d_K|} = m^n\mathfrak{N}(\mathfrak{A})2^{r_2}\sqrt{|d_K|}.$$ 

The result follows. \hfill \Box

In order to prove our goal, we shall need the following:

Lemma 86.8. Let $V$ be a finite dimensional real inner product space and $S \subset V$ be measurable. If $\mathcal{L}$ is a full lattice in $V$ with $\mu(S) > \mu(D_\mathcal{L})$, then there exist $s_1, s_2 \in S$ satisfying $0 \neq s_1 - s_2 \in \mathcal{L}$. 

Proof. The set $S$ is an infinite set, since $\mu(S) > 0$. Suppose that for all $s_1 \neq s_2$ in $S$, we have $s_1 - s_2 \notin \mathcal{L}$. It follows that the the sets
\[(*) \quad S - \lambda, \lambda \in \mathcal{L} \text{ are all disjoint.} \]
Since the measure $\mu$ is translation invariant, for all $\lambda \in \mathcal{L}$, we have $((S - \lambda) \cap D_\mathcal{L}) + \lambda = (D_\mathcal{L} + \lambda) \cap S$. (Check.) It therefore follows that $\mu((S - \lambda) \cap D_\mathcal{L}) = \mu((D_\mathcal{L} + \lambda) \cap S)$. As $V = \bigvee_\mathcal{L}(D_\mathcal{L} + \lambda)$, we have
\[\mu(S) = \sum_{\lambda \in \mathcal{L}} \mu(S \cap (D_\mathcal{L} + \lambda)) = \sum_{\lambda \in \mathcal{L}} ((S - \lambda) \cap D_\mathcal{L}) \leq \mu(D_\mathcal{L})\]
by $(*)$, a contradiction. 

Definition 86.9. Let $V$ be a finite dimensional real inner product space and $X \subset V$. We say:

(i) $X$ is symmetric if $x \in X$, then $-x \in X$.

(ii) $X$ is convex if $x, y \in X$, then $tx + (1 - t)y \in X$ for $0 \leq t \leq 1$, (i.e, the line segment from $x$ to $y$ lies in $X$).

We now establish the goal of this section.

Theorem 86.10. (Minkowski Lattice Point Theorem) Let $V$ be an $n$-dimensional real inner product space and $\mathcal{L}$ be a full lattice in $V$. Suppose that $X \subset V$ is a symmetric, convex, measurable set (respectively, and also compact) and satisfies $\mu(X) > 2^n \mu(D_\mathcal{L})$ (respectively, and $\mu(X) \geq 2^n \mu(D_\mathcal{L})$). Then $X \cap \mathcal{L} \neq \{0\}$.

Proof. Suppose that $\mu(X) > 2^n \mu(D_\mathcal{L})$. Set
\[S := \frac{X}{2} = \{\frac{x}{2} \mid x \in X\}, \text{ so } \mu(S) = \frac{1}{2^n} \mu(X) > \mu(D_\mathcal{L}).\]

By Lemma 86.8, there exist $s_1, s_2 \in S$ such that $0 \neq s_1 - s_2$ lies in $\mathcal{L}$. Let $s_i = \frac{x_i}{2}$ with $\frac{x_i}{2} \in X$ for $i = 1, 2$. So $s_1 - s_2 = \frac{x_1 - x_2}{2}$. By symmetry, $-x_2 \in X$ and by convexity, $s_1 - s_2$, the midpoint of the line segment connecting $x_1$ and $-x_2$, lies in $X$. We are, therefore, done unless $X$ is compact and $\mu(x) = 2^n \mu(D_\mathcal{L})$. Suppose this is the case. Then for each $\epsilon > 0$, the set $(1 + \epsilon)X$ is compact, symmetric, convex, and measurable. Therefore, by what we have seen, there exists a point $0 \neq \lambda_{\epsilon} \in (1 + \epsilon)X \cap \mathcal{L}$. As $(1 + \epsilon)X \cap \mathcal{L}$ is a finite set for all $\epsilon > 0$ by Proposition 86.5, the set $X = \bigcap_{\epsilon>0}(1 + \epsilon)X$ must contain some point $0 \neq \lambda \in X \cap \mathcal{L}$. 

We give a few applications of the Minkowski Lattice Point Theorem.
Application 86.11. Let $f_i = \sum_{j=1}^{n} a_{ij} t_j \in \mathbb{R}[t_1, \ldots, t_n]$, $i = 1, \ldots, n$, be linear forms (i.e., homogeneous of degree one), $\alpha = \det \left( (a_{ij}) \right) \neq 0$, and $c_1, \ldots, c_n > 0$ in $\mathbb{R}$. Suppose that $c_1 \cdots c_n \geq |\alpha|$. Then there exists a nonzero $x \in \mathbb{Z}^n$ satisfying $|f_i(x)| \leq c_i$ for $i = 1, \ldots, n$. (E.g., let $c_i = |\alpha|^{\frac{1}{n}}$.)

**Proof.** Let $\mathcal{L} = \mathbb{Z}^n$ and $X = \{ x \in \mathbb{R}^n \mid |f_i(x)| \leq c_i \}$ for $i = 1, \ldots, n$. Then $X$ is symmetric, convex, compact, and

$$
\mu(X) = \int_{|f_1| \leq c_1} \cdots \int_{|f_n| \leq c_n} dx_1 \cdots dx_n
= \int_{|f_1| \leq c_1} \cdots \int_{|f_n| \leq c_n} \left| \frac{\partial(x_1, \ldots, x_n)}{\partial(f_1, \ldots, f_n)} \right| df_1 \cdots df_n
= \frac{1}{|\alpha|} \int_{|f_1| \leq c_1} \cdots \int_{|f_n| \leq c_n} df_1 \cdots df_n
= \frac{2^n}{|\alpha|} c_1 \cdots c_n \geq 2^n = 2^n \mu(D_L).
$$

The result follows from the Minkowski Lattice Point Theorem. \(\square\)

Recall that the *gamma function* $\Gamma(s)$ for Re$(s) > 0$ is defined by $\Gamma(s) = \int_0^\infty t^{s-1} e^{-t} dt$. We have $\Gamma(s+1) = s\Gamma(s)$, $\Gamma(n+1) = n!$, $\Gamma(\frac{1}{2}) = \sqrt{\pi}$, and $\frac{\pi^{\frac{n}{2}}}{\Gamma(\frac{n}{2}+1)}$ is the volume of the $n$-ball for any integer $n > 0$.

Application 86.12. Let $f_i = \sum_{j=1}^{n} a_{ij} t_j \in \mathbb{R}[t_1, \ldots, t_n]$, $i = 1, \ldots, n$, be linear forms and $\alpha = \det \left( (a_{ij}) \right) \neq 0$. Then there exists a nonzero $x \in \mathbb{Z}^n$ satisfying $f_1^2(x) + \cdots + f_n^2(x) \leq \frac{\alpha}{4} \left( |\alpha| \Gamma(\frac{n}{2}+1) \right)^{\frac{1}{n}}$.

**Proof.** Let $X = \{ x \in \mathbb{R}^n \mid \sum_{i=1}^{n} f_i^2(x) \leq c^2 \}$, a compact, symmetric, convex ellipsoid. Let $\mathcal{L} = \mathbb{Z}^n$ and $c = (\frac{\pi}{\Gamma(\frac{n}{2}+1)} |\alpha|^{\frac{1}{n}})\frac{1}{2^n}$, then

$$
\mu(X) = \int_{\sum_{i=1}^{n} f_i^2 \leq c^2} dx_1 \cdots dx_n = \int_{\sum_{i=1}^{n} f_i^2 \leq c^2} \left| \frac{\partial(x_1, \ldots, x_n)}{\partial(f_1, \ldots, f_n)} \right| df_1 \cdots df_n
= \frac{1}{|\alpha|} \int_{\sum_{i=1}^{n} f_i^2 \leq c^2} df_1 \cdots df_n
= \frac{c^n}{|\alpha|} \left( \frac{\pi^{n/2}}{\Gamma(\frac{n}{2}+1)} \right) = 2^n.
$$

The result follows from the Minkowski Lattice Point Theorem. \(\square\)
Remark 86.13. Let $Q(t_1, \ldots, t_n) = \sum_{i,j=1}^{n} b_{ij} t_i t_j \in \mathbb{R}[t_1, \ldots, t_n]$ be a quadratic form written with $b_{ij} = b_{ji}$ for all $i, j$ (which can always be arranged as $2 \neq 0$). We say $Q$ is **positive definite** if $Q(x) > 0$ for all nonzero $x = (x_1, \ldots, x_n) \in \mathbb{R}^n$. By the Principal Axis Theorem in linear algebra, the matrix $B = (b_{ij}) \in \text{GL}_n(\mathbb{R})$ is diagonalizable with positive eigenvalues (which are squares). In particular, there exists a matrix $A = (a_{ij})$ in $\text{GL}_n(\mathbb{R})$ such that $B = A^t A$. Let $f_i = \sum_{j=1}^{n} a_{ij} t_j$ in $\mathbb{R}[t_1, \ldots, t_n]$, linear forms, and $T^t = (t_1 \cdots t_n)$. Then we have $Q = T^t B T = f_1^2 + \cdots + f_n^2$. The real number $\det B = (\det A)^2$ is called the **discriminant** of $Q$ and written disc $Q$.

So we have:

**Application 86.14.** Let $Q(t_1, \ldots, t_n) = \sum_{i,j=1}^{n} b_{ij} t_i t_j \in \mathbb{R}[t_1, \ldots, t_n]$ be a positive definite quadratic form. Then there exists a nonzero $x \in \mathbb{Z}^n$ satisfying $Q(x) \leq \frac{4}{\pi} \left( \sqrt{\text{disc} Q} \Gamma \left( \frac{n}{2} + 1 \right) \right)^{\frac{2}{n}}$.

We use this to prove a classical theorem of Eisenstein-Hermite. Let $b : \mathbb{R}^n \times \mathbb{R}^n \to \mathbb{R}$ be a symmetric bilinear form, $\mathcal{L} = \mathcal{L}(v_1, \ldots, v_n)$ a full lattice in $\mathbb{R}^n$ satisfying $b|_{\mathcal{L} \times \mathcal{L}} : \mathcal{L} \times \mathcal{L} \to \mathbb{Z}$ with $B_b := ((b(v_i, v_j))$ in $\text{GL}_n(\mathbb{R})$, i.e., $\det B_b \neq 0$. We then say that $(\mathcal{L}, b)$ is a **$\mathbb{Z}$-bilinear space**. We call such a bilinear space **positive definite** if $Q_b = T^t B_b T$ is positive definite. This is equivalent to $q_b : V \to \mathbb{R}$ defined by $q_b(x) = b(x, x) > 0$ for all $0 \neq x \in \mathbb{R}^n$. We call $n$ the **rank** of $\mathcal{L}$ and $\det B_b$ the **discriminant** of $(\mathcal{L}, b)$. It is unique up to a $\mathbb{Z}$-isometry (i.e., a $\mathbb{Z}$-isomorphism preserving $b$) as a change of basis matrix has determinant $\pm 1$.

Note that if $(\mathcal{L}, b)$ is a positive definite $\mathbb{Z}$-bilinear space, there exists a nonzero $x \in \mathbb{Z}^n$ satisfying Application 86.14.

**Theorem 86.15.** (Eisenstein-Hermite) For each pair of integers $(n, d)$ with $n > 0$, there exist finitely many positive definite $\mathbb{Z}$-bilinear spaces of rank $n$ and discriminant $d$ up to isometry.

**Proof.** Let $(\mathcal{L}, b)$ be a positive definite $\mathbb{Z}$-bilinear space of rank $n$ and discriminant $d = \text{disc} \mathcal{L}$. Set $c(n, d) = \frac{4}{\pi} \left( \Gamma \left( \frac{n}{2} + 1 \right) \right)^{\frac{2}{n}} \sqrt{d}$. By Application 86.14, there exists a nonzero $x \in \mathcal{L}$ that satisfies $0 < q(x) := b(x, x) \leq c(n, d)$. Let

$$\mathcal{L}_0 := \{ y \in \mathcal{L} \mid b(x, y) \equiv 0 \mod q(x) \} \subset \mathcal{L}.$$ 

Clearly, we have $[\mathcal{L} : \mathcal{L}_0] \leq q(x)$.

**Claim.** $\mathcal{L}_0 = \mathbb{Z}x \oplus (\mathbb{Z}x)^\perp$ over $\mathbb{Z}$ (with the obvious definition):
As \((\mathcal{L}, b)\) is positive definite, we have \(\mathbb{Z} x \cap (\mathbb{Z} x)^\perp = 0\). Let \(y \in \mathcal{L}_0\). By definition, \(b(x, y) = a q(x)\) for some \(a \in \mathbb{Z}\) and \(b(x, y - ax) = 0\). Consequently, \(y = (y - ax) + ax\) lies in \((\mathbb{Z} x)^\perp + \mathbb{Z} x\) proving the claim.

In the notation of Remark 86.13, we had \(B_b = A^t A\). This translates into

\[
\text{disc } \mathcal{L}_0 = \mu(D_{\mathcal{L}_0})^2 = [\mathcal{L} : \mathcal{L}_0]^2 \mu(D_{\mathcal{L}})^2 = [\mathcal{L} : \mathcal{L}_0]^2 \text{disc } \mathcal{L}
\]

using Lemma 86.6. It follows that

\[
\text{disc } \mathbb{Z} x \text{ disc}(\mathbb{Z} x)^\perp = \text{disc } \mathcal{L}_0 \leq q(x)^2 d \leq c(n, d)^2 d.
\]

By induction on rank \(\mathcal{L}\), there exist finitely many positive definite \(\mathbb{Z}\)-bilinear spaces \((\mathbb{Z} x)^\perp, \mathbb{Z} x\) (up to isometry) of discriminant all bounded by \(c(n, d)^2 d\). Hence there exist finitely many positive definite \(\mathbb{Z}\)-bilinear spaces \(\mathcal{L}_0\) (up to isometry). Since \(\mathcal{L}/\mathcal{L}_0\) is a finite group and \(b\) is determined by its values on a basis for \(\mathcal{L}\), there exist only finitely many possible \((\mathcal{L}, b)\) (up to isometry). \(\square\)

We shall see in Corollary 88.14 below that there exist finitely many number fields \(K\) with fixed discriminant \(d\).

**Exercises 86.16.**

1. Use the Minkowski Lattice Point Theorem to prove Fermat’s Theorem that positive primes congruent to one modulo four are sums of two squares.

2. Use the Minkowski Lattice Point Theorem to prove Lagrange’s Theorem that every positive integer is a sum of four squares.

3. Prove an analogous result as in the Eisenstein-Hermite Theorem for indefinite \(\mathbb{Z}\)-bilinear spaces, i.e., those \(\mathbb{Z}\)-bilinear spaces \((\mathcal{L}, b)\) in which there exist \(x, y \in \mathcal{L}\) satisfying \(b(x, x) < 0\) and \(b(y, y) > 0\).

4. Let \(f_i = \sum_{j=1}^{n} a_{ij} t_j \in \mathbb{R}[t_1, \ldots, t_n], i = 1, \ldots, n\), be linear forms, \(\alpha = \det ((a_{ij})) \neq 0\), and \(r \in \mathbb{R}\). Show there exists a nonzero \(x \in \mathbb{Z}^n\) satisfying \(\sum_{i=1}^{n} |f_i(x)| \leq r\) if \(r \geq (n! |\alpha|)^{\frac{1}{n}}\).

### 87. Units in a Ring of Algebraic Numbers

In Section 83, we computed the units in the ring of algebraic numbers in a quadratic number field \(\mathbb{Q}(\sqrt{m})\), \(m\) a square-free integer. In this section, we generalize this result by proving the Dirichlet Unit Theorem that determines the units in \(\mathbb{Z}_K\) for any algebraic number field \(K\). We continue to use the notation of Section 86.

We begin with a computation.
Lemma 87.1. (Minkowski) Let $A = (a_{ij}) \in \mathbb{M}_n(\mathbb{R})$ satisfy $a_{ij} < 0$ for all $i, j = 1, \ldots, n$ and $i \neq j$. Suppose that $\sum_{j=1}^{n} a_{ij} > 0$ for each $i = 1, \ldots, n$. Then $A \in \text{GL}_n(\mathbb{R})$.

Proof. (Artin). Suppose that $\det A = 0$. Let $v_1, \ldots, v_n$ denote the columns of $A$. Then $v_1, \ldots, v_n$ are linearly dependent, so there exist $x_1, \ldots, x_n \in \mathbb{R}$ not all zero such that $\sum_{i=1}^{n} x_i v_i = 0$, i.e., $\sum_{i=1}^{n} x_i a_{ii} = 0$ for $l = 1, \ldots, n$. Choose $l$ such that $|x_l|$ is maximal. Multiplying by $-1$ if necessary, we may assume that $x_l > 0$. Since $x_l \geq x_i$ and $a_{ii} < 0$ for all $i \neq l$, $i = 1, \ldots, n$, we have

$$0 = x_l a_{ll} + \sum_{i=1 \atop i \neq l}^{n} x_i a_{il} \geq x_l a_{ll} + \sum_{i=1 \atop i \neq l}^{n} x_l a_{ii} = x_l \sum_{i=1}^{n} a_{ii} > 0,$$

a contradiction. □

Construction/Definition 87.2. Let $K$ be a number field of degree $n$. As we have $n$ embeddings of $K$ with $\sigma_{r_1+r_2+j} = \sigma_{r_1+j}$ for $j = 1, \ldots, r_2$, we set

$$e_i := \begin{cases} 1 & \text{if } 1 \leq i \leq r_1 \\ 2 & \text{if } 1 \leq r_1 + 1 \leq i \leq r_1 + r_2 \end{cases}$$

(with $r_1$ or $r_2 = 0$ possible), as we shall be interested in norms from $K$.
If $x \in K$, define

$$\ell_i(x) = e_i \log |x^{(i)}|, \quad 1 \leq i \leq r_1 + r_2$$

and the logarithmic map

$$\ell : K^x \to \mathbb{R}^{r_1+r_2} \quad \text{by} \quad x \mapsto (\ell_1(x), \ldots, \ell_{r_1+r_2}(x)).$$

In particular, setting

$$U_K := \mathbb{Z}_K^x, \quad \text{we have if } \varepsilon \in U_K, \text{ then } N_{K/\mathbb{Q}}(\varepsilon) = \prod_{i=1}^{n} \varepsilon^{(i)} = \pm 1.$$ 

Taking $\log(\prod_{i=1}^{n} |\varepsilon^{(i)}|)$, we get

$$(87.3) \quad 0 = \sum_{i=1}^{r_1+r_2} e_i \log |\varepsilon^{(i)}|.$$ 

It follows that $\{\log |\varepsilon^{(1)}|, \ldots, \log |\varepsilon^{(r_1+r_2)}|\}$ are linearly dependent. We eliminate $\log |\varepsilon^{(r_1+r_2)}|$ from this linearly dependent set. i.e., $\varepsilon^{(r_1+r_2)}$ and set $r = r_1 + r_2 - 1$. So $\ell(U_K) \subset \mathbb{R}^r$ lying on the hyperspace in $\mathbb{R}^{r_1+r_2}$ defined by $0 = \sum_{i=1}^{r_1+r_2} t_i$ in $\mathbb{R}^{r+1}$. 
Lemma 87.4. Let $K$ be a number field of degree $n$ and $r = r_1 + r_2 - 1$. Suppose that $C \subset \mathbb{R}^{r_1 + r_2}$ is a compact set. Then $|\ell^{-1}(C \cap \ell(U_K))| < \infty$. In particular, $\ell(U_K) \subset \mathbb{R}^r$ is a lattice (although not a priori full).

Proof. We have $C$ is compact and equation (87.3) holds. It follows that if $\ell(\varepsilon) \in C$ for $\varepsilon \in U_K$ then all the coordinates of such an $\ell(\varepsilon)$ are uniformly bounded, i.e., there exists an $M_0 > 0$ such that $\ell_i(\varepsilon) \leq M_0$ for $i = 1, \ldots, r_1 + r_2$ for any $\varepsilon \in U_K$ satisfying $\ell(\varepsilon) \in C$. In particular, $|\varepsilon^{(i)}| < e^{M_0}$, for $i = 1, \ldots, r_1 + r_2$. As $M_1 = e^{M_0} > 1$, we have $|\varepsilon^{(i)}| < M_1$, for $i = 1, \ldots, n$. Let $Y = \{x \in \mathbb{R}^{r_1} \times \mathbb{C}^{r_2} | |x_i| \leq M_1\}$, a compact subset of $\mathbb{R}^{r_1} \times \mathbb{C}^{r_2}$ containing $i_C(\ell^{-1}(C \cap \ell(U_K)))$. Since $i_C(\mathbb{Z}_K)$ is a lattice in $\mathbb{R}^{r_1} \times \mathbb{C}^{r_2}$, we have $|i_C(\ell^{-1}(C \cap \ell(U_K)))| < \infty$ by Proposition 86.5. Hence $\ell(U_K) \subset \mathbb{R}^r$ is a lattice by Proposition 86.5.

Corollary 87.5. Let $K$ be a number field. Then $\ker \ell$ is a finite cyclic group.

Proof. Any finite subgroup of $K^\times$ is cyclic.

Notation 87.6. Let $K$ be a number field. Let $W_K := \ker \ell$, and (as usual) let $\mu_K$ denote the roots of unity in $K$.

We can now classify the units in a ring of algebraic integers.

Theorem 87.7. (Dirichlet Unit Theorem) Let $K$ be a number field of degree $n = r_1 + 2r_2$ and $r = r_1 + r_2 - 1$. Let $\ell : U_K \to \mathbb{R}^r$ be the logarithmic map given by $\varepsilon \mapsto (\ell_1(\varepsilon), \ldots, \ell_r(\varepsilon))$. Then

1. $W_K = \mu_K$ is a finite cyclic group.
2. $\ell(U_K) \subset \mathbb{R}^r$ is a full lattice.
3. $U_K \cong W_K \times \mathbb{Z}^r$.

Proof. (1): If $w \in W_K$, then $|w|_{W_K} = 1$ by Corollary 87.5, so $W_K \subset \mu_K$. Conversely if $x \in U_K$ satisfies $x^m = 1$ for some $m > 0$, then $\ell_i(x) = 0$ for $i = 1, \ldots, n$, hence $x \in W_K$.

(2): We know that $\ell(U_K) \subset \mathbb{R}^r$ is a lattice by Lemma 87.4, so we need to show that it is a full lattice. To do this it suffices to find $\varepsilon_1, \ldots, \varepsilon_r$ in $U_K$ such that $\mathcal{L}(\ell(\varepsilon_1), \ldots, \ell(\varepsilon_r))$ is a full lattice in $\mathbb{R}^r$.

Let $V = \mathbb{R}^{r_1} \times \mathbb{C}^{r_2}$ and define $N : V \to \mathbb{R}$ by

$$N(x) = x_1 \cdots x_{r_1}x_{r_1+1} \cdots x_{r_1+r_2}x_{r_1+r_2+1} \cdots x_{r_1+r_2}.$$

Viewing $i_C : K \hookrightarrow V$ as an inclusion, we have $N|_K = N_{K/Q}$. Let $\{w_1, \ldots, w_n\}$ be an integral basis for $\mathbb{Z}_K$. Therefore, $\mathbb{Z}_K \hookrightarrow V$ is a full lattice in $V$ viewing $w_j = (w_j^{(1)}, \ldots, w_j^{(r_1+r_2)})$. In particular, $V$ is
a $\mathbb{Z}_K$-module via $\alpha(v_1, \ldots, v_{r_1+r_2}) = (\alpha^{(1)}v_1, \ldots, \alpha^{(r_1+r_2)}v_{r_1+r_2})$ for all $\alpha \in \mathbb{Z}_K$. It follows if $x \in V$ satisfies $N(x) \neq 0$, then $\mathcal{L}_x := \mathbb{Z}_K x = \mathbb{Z}w_1 x \oplus \cdots \oplus \mathbb{Z}w_n x$ is a full lattice in $V$. Using an analogous calculation that we did to prove Proposition 86.7, we see that if $N(x) \neq 0$, then

$$\text{Vol}(V/\mathcal{L}_x) = |N(x)| 2^{-r_2} \sqrt{|d_K|}.$$ 

In particular, if $N x = 1$, then

$$\text{Vol}(V/\mathcal{L}_x) = 2^{-r_2} \sqrt{|d_K|}$$

is independent of $x \in V$ satisfying $N(x) = 1$.

**Claim.** Let $X$ be a compact, convex, symmetric, measurable subset of $V$ with $\mu(X) > 2^{r_2} \sqrt{d_K}$. Then

1. There exists $M_0 = M_0(X) > 0$ such that if $\beta \in X$, $|\beta_i| < M_0$, $i = 1, \ldots, n$.

2. If $x \in V$ satisfying $N(x) = 1$, then there exists $0 \neq \alpha \in \mathbb{Z}_K$ such that $\alpha x \in \mathcal{L}_x \cap X$ and an $M = M(X) > 0$ independent of $x$ satisfying $|N_{K/Q}(\alpha)| < M$.

The first statement is immediate as $X$ is compact. For the second, fix $x \in V$ that satisfies $N(x) = 1$. By the Minkowski Lattice Point Theorem, there exists $0 \neq \alpha \in \mathbb{Z}_K$ satisfying $\alpha x \in \mathcal{L}_x \cap X$. As $X$ is compact, by (1) all the coordinates of $\alpha x$ are uniformly bounded by a constant $M_0 = M_0(X) > 0$. Since $N(\alpha x) = N_{K/Q}(\alpha) N(x) = N_{K/Q}(\alpha)$, $M = M_0^n$ works and the claim is proven.

By Proposition 85.8, we have

$$N_{K/Q}(\alpha) \mathbb{Z} = N_{K/Q}(\alpha \mathbb{Z}_K) = \prod_p N_{K/Q}(p)^{\nu_p(\alpha)} = \prod_p \mathcal{N}(p)^{\nu_p(\alpha)} \mathbb{Z}.$$ 

Consequently, if $0 \neq \alpha \in \mathbb{Z}_K$ satisfies $|N_{K/Q}(\alpha)| < M$, then $\mathcal{N}(p) < M$ for all prime ideals $p$ such that $v_p(\alpha) > 0$ and the allowable values $v_p(\alpha) > 0$ are also bounded. It follows that there exist finitely many principal integral ideals $\alpha \mathbb{Z}_K$ (so $\alpha \in \mathbb{Z}_K$) satisfying $0 < |N_{K/Q}(\alpha)| < M$.

Let $\alpha_1 \mathbb{Z}_K, \ldots, \alpha_s \mathbb{Z}_K$

be all the finitely many nonzero principal integral ideals that satisfy $|N_{K/Q}(\alpha_i)| < M$ for $i = 1, \ldots, s$.

[Note that if $|U_K|$ is infinite, then $|N_{K/Q}(\alpha u)| = |N_{K/Q}(\alpha)|$ for any algebraic integer $\alpha$ and any $u \in U_K$. In particular, in the above, it is crucial that we deal with principal ideals rather than elements.]

Let $\alpha$ be any element that satisfies the claim. Then there exists $k$, $1 \leq k \leq s$, such that $\alpha \mathbb{Z}_K = \alpha_k \mathbb{Z}_K$ and an $\varepsilon \in U_K$ with $\alpha = \varepsilon \alpha_k$. 

Therefore, for each $x$ satisfying $N(x) = 1$, we have:

(87.8)

There exists an $\varepsilon \in U_K$ and a $k$, $1 \leq k \leq s$, with $\varepsilon^k x \in X$.

In particular, $|\varepsilon^{(i)} \alpha_0^{(i)} x_i| < M_0$, $i = 1, \cdots, r_1 + r_2$.

Fix $j$ for $1 \leq j \leq r_1 + r_2$ and choose $x \in V$ with $N(x) = 1$ and $|x_i| >> 0$ for all $i \neq j$, $i = 1, \ldots, r_1 + r_2$.

Let $\varepsilon_j$ be the corresponding unit in (87.8).

Since there exist finitely many $\alpha_0^{(i)}$, in (87.8), we may assume that we have chosen $x$ with the $|x_i|$, $i \neq j$ so large, $1 \leq i \leq r_1 + r_2$, that 87.8 also implies

(*) $|\varepsilon_j^{(i)}| < 1$ i.e., $\ell_i(\varepsilon_j) < 0$ for all $1 \leq i \leq r_1 + r_2$.

For each $j \leq r$, choose such $x \in V$ and corresponding $\varepsilon_j$ satisfying (87.8) and (*). Set

$$
\beta = \begin{pmatrix}
\ell_1(\varepsilon_1) & \cdots & \ell_r(\varepsilon_1) \\
\vdots & \ddots & \vdots \\
\ell_1(\varepsilon_r) & \cdots & \ell_r(\varepsilon_r)
\end{pmatrix}.
$$

As $\ell_i(\varepsilon_j) < 0$ for $i \neq j$ and

$$
0 = \sum_{i=1}^{r_1 + r_2} \ell_i(\varepsilon_j) = \sum_{i=1}^{r} \ell_i(\varepsilon_j) + \ell_{r+1}(\varepsilon_j),
$$

we have for all $j \leq r$,

$$
\sum_{i=1}^{r} \ell_i(\varepsilon_j) = -\ell_{r+1}(\varepsilon_j).
$$

By Minkowski’s Lemma 87.1, $\det \beta \neq 0$. Therefore, $\ell(\varepsilon_1), \ldots, \ell(\varepsilon_r)$ are linearly independent as needed.

(3): Since $\ell(U_K)$ is $\mathbb{Z}$-free, (3) follows from the split exact sequence

$$
1 \to W_K \to U_K \to \ell(U_K) \to 1.
$$

**Remark 87.9.** Let $K$ be a number field. The proof of the Dirichlet Unit Theorem shows how to get $r$ linearly independent units generating a subgroup of finite index in $U_K$. To get all units, one must use some type of descent. This is usually not easy.

**Definition 87.10.** Let $K$ be a number field of degree $n$ over $\mathbb{Q}$ and $\{\varepsilon_1, \ldots, \varepsilon_r\}$ a basis for the free part of $U_K$. We call this basis a set of
base units. Therefore, if \( \nu \in U_K \), there exist unique \( \nu_i \geq 0 \), \( i = 1, \ldots, r \), and a unit \( w \in W_K \) such that \( \nu = w\varepsilon_1 \cdots \varepsilon_r \). We let
\[
R_K := |\det (\ell_i(\varepsilon_j))|, \quad \text{for } 1 \leq i, j \leq r
\]
called the regulator of \( K \).

Remark 87.11. Let \( K \) be a number field of degree \( n \) over \( \mathbb{Q} \).

1. \( R_K \) is independent of the choice of the last embedding that we drop.

**Proof.** Let \( A = 
\begin{pmatrix}
\ell_1(\varepsilon_1) & \cdots & \ell_r(\varepsilon_1) & \ell_{r+1}(\varepsilon_1)
\vdots & & \vdots & \\
\ell_1(\varepsilon_r) & \cdots & \ell_r(\varepsilon_r) & \ell_{r+1}(\varepsilon_r)
\end{pmatrix}
\begin{pmatrix}
e_1 \\
e_r \\
e_1 \\
0
\end{pmatrix}
\]
with
\[
e_i := \begin{cases} 
1 & \text{if } 1 \leq i \leq r_1 \\
2 & \text{if } 1 \leq r_1 < i \leq r_1 + r_2
\end{cases}
\]
as usual. Then \( A \) and the matrices
\[
\begin{pmatrix}
\ell_1(\varepsilon_1) & \cdots & \ell_r(\varepsilon_1) & \sum_{i=1}^{r+1} \ell_i(\varepsilon_1)
\vdots & & \vdots & \\
\ell_1(\varepsilon_r) & \cdots & \ell_r(\varepsilon_r) & \sum_{i=1}^{r+1} \ell_i(\varepsilon_r)
\end{pmatrix}
\begin{pmatrix}
e_1 \\
e_r \\
e_1 \\
0
\end{pmatrix}
\text{ and }
\begin{pmatrix}
\ell_1(\varepsilon_1) & \cdots & \ell_r(\varepsilon_1) & 0 
\vdots & & \vdots & \\
\ell_1(\varepsilon_r) & \cdots & \ell_r(\varepsilon_r) & 0 
\end{pmatrix}
\begin{pmatrix}
e_1 \\
e_r \\
e_1 \\
n
\end{pmatrix}
\]
all have the same determinant. Hence \( \det A = (-1)^{n-1}nR_K \). It follows that \( R_K \) is independent of the coordinate that is dropped. \( \square \)

2. \( R_K \) is independent of the choice of base units as a change of basis matrix has determinant equal to \( \pm 1 \).
3. \( R_K \) is very important but an elusive invariant.

We next indicate how the Dirichlet Unit Theorem generalizes. We leave the proofs as exercises.

Definition 87.12. Let \( K \) be a number field of degree \( n \) over \( \mathbb{Q} \). We call the set of nonzero prime ideals, \( \text{Max}(\mathbb{Z}_K) \), in \( \mathbb{Z}_K \), the set of finite primes in \( K \) and the set \( S_\infty := \{\sigma_1, \ldots, \sigma_n\} \) of embeddings \( K \to \mathbb{C} \) the set of infinite primes in \( K \) with \( \{\sigma_1, \ldots, \sigma_{r_1}\} \) the set of real infinite primes and \( \{\sigma_{r_1+1}, \ldots, \sigma_n\} \) the set of complex infinite primes. We let \( S_K = \text{Max}(\mathbb{Z}_K) \cup S_\infty(K) \) called the set of extended primes of \( \mathbb{Z}_K \). [These represent all the necessary completions of \( K \) – the finite primes
88. Minkowski Bound

In this section, we use the Minkowski Lattice Theorem in \( \mathbb{R}^{r_1} \times \mathbb{C}^{r_2} \) to show that the class number \( h_K := |I_{\mathbb{Z}_K}/P_{\mathbb{Z}_K}| \) of a ring of algebraic integers in the number field \( K \) is finite. We shall also show that there exist only finitely many algebraic number fields with a given discriminant. We continue to use the notation of Section 86.

We investigate the following subset of \( V = \mathbb{R}^{r_1} \times \mathbb{C}^{r_2} \): If \( x \in V \), then \( x = (x_1, \ldots, x_{r_1}, z_{r_1+1}, \ldots, z_{r_1+r_2}) \). Set

\[
X(t) = X_{r_1, r_2}(t) := \{ x \in V \mid \sum_{i=1}^{r_1} |x_i| + 2 \sum_{i=r_1+1}^{r_1+r_2} |z_i| \leq t \}
\]

\[
= \{ x \in V \mid \sum_{i=1}^{r_1} |x_i| + 2 \sum_{i=1}^{r_2} \sqrt{x_{r_1+i}^2 + y_{r_1+i}^2} \leq t \}.
\]

[Note this is the same as

\[
X(t) := \{ x \in V \mid \sum_{i=1}^{n} |x_i| \leq t, \ x_{r_1+r_2+j} = \pi_{r_1+j} \text{ for } j = 1, \ldots, r_2 \}.
\]

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\]

\[
= \{ x \in V \mid \sum_{i=1}^{r_1} |x_i| + 2 \sum_{i=1}^{r_2} \sqrt{x_{r_1+i}^2 + y_{r_1+i}^2} \leq t \}.
\]

[Note this is the same as

\[
X(t) := \{ x \in V \mid \sum_{i=1}^{n} |x_i| \leq t, \ x_{r_1+r_2+j} = \pi_{r_1+j} \text{ for } j = 1, \ldots, r_2 \}.
\]
Lemma 88.1. $X(t)$ is a compact, convex, symmetric set with measure

$$\mu(X(t)) = \frac{2^{r_1}}{n!} \left(\frac{\pi}{2}\right)^{r_2} t^n.$$ 

Proof. Clearly, $X(t)$ is compact and symmetric. Let $0 \leq a \leq 1$, $b = 1 - a$ and $x = (x_1, \ldots, x_n)$, $y = (y_1, \ldots, y_n)$ in $X(t)$. Then $\sum_{i=1}^n |ax_i + by_i| \leq \sum_{i=1}^n (a|x_i| + b|y_i|) \leq at + bt = t$. Therefore, $X(t)$ is convex. It remains to determine $\mu(X(t))$.

Set $X_0(t) = 1$. Changing variables shows that

($\ast$) \quad $\mu(X_{r_1,r_2}(t)) = t^n \mu(X_{r_1,r_2}(1))$,

so we need only compute $\mu(X_{r_1,r_2}(1))$. Suppose that $r_1 > 0$. Then by ($\ast$) and induction, we have

$$\mu(X_{r_1,r_2}(1)) = 2 \int_0^1 \mu(X_{r_1-1,r_2}(1-t)) \, dt$$

$$= 2 \mu(X_{r_1-1,r_2}(1)) \int_0^1 (1-t)^{r_1-1+2r_2} \, dt$$

$$= \frac{2^{r_1}}{n} \mu(X_{r_1-1,r_2}(1)) = \frac{2^{r_1}}{(r_1 + 2r_2) \cdots (2r_2 + 1)} \mu(X_{r_1,r_2}(1)).$$

Therefore, it suffices to compute $\mu(X_{r_1,r_2}(1))$. We have

$$\mu(X_{r_1,r_2}(1)) = \iint_{x^2 + y^2 \leq \frac{1}{4}} \mu(X_{r_1,r_2-1}(1 - 2\sqrt{x^2 + y^2})) \, dxdy.$$  

Let $\rho = \sqrt{x^2 + y^2}$, $x = \rho \cos \theta$ be polar coordinates. Then by ($\ast$) and letting $u = 1 - 2\rho$, we have

$$\mu(X_{r_1,r_2}(1)) = \int_{\theta=0}^{\theta=2\pi} \int_{\rho=0}^{\rho=\frac{1}{2}} \mu(X_{r_1,r_2-1}(1 - 2\rho)) \, \rho \, d\rho$$

$$= 2\pi \mu(X_{r_1,r_2-1}(1)) \int_{\rho=0}^{\rho=\frac{1}{2}} (1 - 2\rho)^{2(r_2-1)} \, d\rho$$

$$= \frac{\pi}{2} \mu(X_{r_1,r_2-1}(1)) \int_{u=1}^{u=0} u^{2(r_2-1)} (1 - u) \, du$$

$$= \frac{\pi}{2} \mu(X_{r_1,r_2-1}(1)) \frac{1}{2r_2(2r_2 - 1)} = \left(\frac{\pi}{2}\right)^{r_2} \frac{1}{(2r_2)!},$$

and we are done. \qed

Definition 88.2. Define the Minkowski Bound to be

$$B_{n,r_2} := \left(\frac{4}{\pi}\right)^{r_2} \frac{n!}{n^n}.$$
If \( K \) is a number field of degree \( n_K \), we let \( B_K = B_{n_K, r_2^K} \) where \( n_K = r_1^K + 2r_2^K \).

Note that \( B_K \to 0 \) as \( n_K \to \infty \).

The key result in this section is the following crucial inequality.

**Theorem 88.3.** Let \( K \) be a number field and \( \mathfrak{a} \in I_{\mathbb{Z}_K} \). Then there exists a nonzero \( a \in \mathfrak{a} \) satisfying:

\[
|N_{K/\mathbb{Q}}(a)| \leq B_K \mathfrak{m}(\mathfrak{a}) \sqrt{|d_K|}.
\]

**Proof.** Let \( X(t) = X_{r_1, r_2}(t) \subset V = \mathbb{R}^{r_1} \times \mathbb{C}^{r_2} \) be as above. Choose \( t \) to satisfy \( \mu(X(t)) = 2^n \mu(V/\mathfrak{a}) \). By Proposition 86.7, we then have

\[
\frac{2^{r_1}}{n!^2} \frac{\pi}{2}^{r_2} t^n = \mu(X(t)) = 2^n \mu(V/\mathfrak{a})
= 2^n \mathfrak{m}(\mathfrak{a}) 2^{-r_2} \sqrt{|d_K|} = 2^{r_1 + r_2} \mathfrak{m}(\mathfrak{a}) \sqrt{|d_K|}.
\]

Therefore,

\[
t^n = \left( \frac{4}{\pi} \right)^{r_2} \frac{n!^2}{2} \mathfrak{m}(\mathfrak{a}) \sqrt{|d_K|}.
\]

By the Minkowski Lattice Point Theorem, there exists a nonzero \( \beta \) in \( \mathfrak{a} \cap X(t) \). Therefore, \( \sum_{i=1}^n |\beta(i)| \leq t \). Since the geometric mean is bounded by the arithmetic mean, we have \( \left( \prod_{i=1}^n |\beta(i)| \right)^{\frac{1}{n}} \leq \frac{1}{n} \sum_{i=1}^n |\beta(i)| \).

Hence \( |N_{K/\mathbb{Q}}(\beta)| \leq \frac{t^n}{n^n} \leq B_K \mathfrak{m}(\mathfrak{a}) \sqrt{|d_K|} \) as needed. \( \square \)

This theorem has important applications.

**Definition 88.4.** Let \( K \) be a number field. Recall the **class group** of \( \mathbb{Z}_K \) is \( Cl_{\mathbb{Z}_K} := I_{\mathbb{Z}_K}/P_{\mathbb{Z}_K} \). We let \( h_K := |Cl_{\mathbb{Z}_K}| \), called the **class number** of \( K \) (or \( \mathbb{Z}_K \)).

**Corollary 88.5.** Let \( K \) be a number field and \( \mathfrak{a} \in Cl_{\mathbb{Z}_K} \). Then there exists a nonzero integral ideal \( \mathfrak{a} \subset \mathbb{Z}_K \) in \( \mathfrak{a} \) satisfying \( \mathfrak{m}(\mathfrak{a}) \leq B_K \sqrt{|d_K|} \).

**Proof.** Let \( \mathfrak{b} \in \mathfrak{a} \). Then by the theorem, there exists a nonzero \( a \in \mathfrak{b}^{-1} \) satisfying \( |N_{K/\mathbb{Q}}(a)| \leq B_K \mathfrak{m}(\mathfrak{b}^{-1}) \sqrt{|d_K|} \). Since \( \mathfrak{m}(a\mathbb{Z}_K) = |N_{K/\mathbb{Q}}(a)| \) by Corollary 85.9, we have \( \mathfrak{m}(a\mathfrak{b}) \leq B_K \sqrt{|d_K|} \). As \( a \in \mathfrak{b}^{-1} \), the fractional ideal \( a\mathfrak{b} \subset \mathbb{Z}_K \) is an ideal. It follows that \( \mathfrak{a} = a\mathfrak{b} \in \mathfrak{a} \) works. \( \square \)

**Corollary 88.6.** Let \( K \) be a number field. Then the class number \( h_K \) is finite.
Proof. The inequality $\mathfrak{N}(\mathfrak{A}) \leq B_K \sqrt{|d_K|}$ is satisfied by finitely many ideals $\mathfrak{A} \subset \mathcal{O}_K$, since only finitely many ideals can contract to any fixed $m\mathbb{Z}$ by the unique factorization of ideals in a Dedekind domain. □

Remark 88.7. The class group of a Dedekind domain may not be finite. For example, the Dedekind domain $E = \mathbb{C}[t, t_1]/(t_1^2 - t(t-1)(t+1))$, the integral closure of $\mathbb{C}[t]$ in the quotient field of $E$, corresponding to the points on the elliptic curve $y^2 = x(x-1)(x+1)$ has an infinite class group.

Corollary 88.8. Let $K$ be a number field. Suppose that for every positive prime integer $p$ satisfying $p \leq B_K \sqrt{|d_K|}$, any prime $p$ in $\mathcal{O}_K$ is principal if it satisfies $p\mathbb{Z} = p \cap \mathbb{Z}$. Then $h_K = 1$. In particular, $\mathcal{O}_K$ is a UFD.

Proof. Let $A \in \text{Cl}_{\mathcal{O}_K}$. Choose $\mathfrak{A} \in A$ an integral ideal satisfying $\mathfrak{N}(\mathfrak{A}) \leq B_K \sqrt{|d_K|}$. Let $\mathfrak{N}(\mathfrak{A}) = p_1 \cdots p_n$ be a factorization into primes in $\mathcal{O}_K$. As $\mathfrak{N}$ is multiplicative, the hypothesis implies that $\mathfrak{A}$ is principal. So $\text{Cl}_{\mathcal{O}_K} = 1$. □

Example 88.9. Let $K = \mathbb{Q}(\sqrt{m})$ with $m$ a square-free integer. Then $h_K = 1$ if $m = -1, -2, -3, -7, -11, -19, -43, -67, -163$.

Baker and Stark showed that the fields in the example are the only imaginary quadratic fields $K$ with $h_K = 1$.

Corollary 88.10. Let $K$ be a number field of degree greater than one. Then $|d_K| > 1$. In particular, $\mathcal{O}_K$ is not unramified over $\mathbb{Z}$.

Proof. We know that $\left( \frac{\pi}{4} \right)^2 \frac{n^n}{n!} > 1$ if $n > 1$ (as the left hand side is increasing). Therefore, there exists a nonzero integral ideal $\mathfrak{A} \subset \mathcal{O}_K$ satisfying

$$\sqrt{|d_K|} \geq \frac{\mathfrak{N}(\mathfrak{A})}{B_K} \geq \left( \frac{\pi}{4} \right)^2 \frac{n^n}{n!} = \left( \frac{\pi}{4} \right)^2 \frac{n^n}{n!} > 1.$$ (88.11)

Therefore, $K$ ramifies over $\mathbb{Z}$ by the Dedekind Ramification Theorem 82.2. □

This result is special to $\mathbb{Q}$. If we look at relative number theory, i.e., $L/K/\mathbb{Q}$ is finite with $K \neq \mathbb{Q}$, it can happen that $\mathcal{O}_L$ is unramified over $\mathcal{O}_K$ with $K < L$. For any such field $K$, there exists a unique number field $L$ with $L/K$ the maximal abelian extension of $K$ that satisfies $\mathcal{O}_L$ is unramified over $\mathcal{O}_K$. It is called the Hilbert class field of $K$. It has the properties that its degree over $K$ is $h_K$ and $G(L/K) \cong \text{I}_{\mathcal{O}_K}$ (canonically). So Corollary 88.10 says that $\mathbb{Q}$ is Hilbert class field of $\mathbb{Q}$. 
For the next consequence, we need Stirling’s Formula (that we do not prove).

**Fact 88.12.** (Stirling’s Formula) Let $N >> 0$ be an integer. Then

$$\sqrt{2\pi N} \ N^N e^{-N} < N! < \sqrt{2\pi N} \ N^N e^{-N} \left(1 + \frac{1}{12N-1}\right).$$

**Corollary 88.13.** Let $K$ be a number field of degree $n_K > 1$. Then $n_K \log(|d_K|)$ is bounded.

**Proof.** By equation (88.11) (and squaring), we have, with $n = n_K$,

$$\frac{n}{\log(|d_K|)} \leq \frac{n}{2 \log \left(\frac{n\pi}{4} \left(\frac{\pi}{4}\right)^\frac{n}{2}\right)}.$$

Therefore, by Stirling Formula, we have

$$\frac{n}{\log(|d_K|)} \leq \frac{n}{2 \log \left(\frac{e^n}{\sqrt{2\pi n}} \left(\frac{\pi}{4}\right)^\frac{n}{2} \left(1 + \frac{1}{12n-1}\right)\right)}.$$

$$\leq \frac{2n - \log(2\pi n) + n \log\left(\frac{n}{4}\right) + \log\left(1 + \frac{1}{12n-1}\right)}{n}$$

which is bounded (say by L’Hôpital’s Rule).

**Corollary 88.14.** There exists finitely many number fields $K$ with discriminant $d$.

**Proof.** Fix a field $K$ and set $n = n_K$. We use the embedding $i_C : K \to \mathbb{R}^{r_1} \times \mathbb{C}^{r_2}$. 

**Case 1.** $r_2 > 0$:

Define $X \subset \mathbb{R}^{r_1} \times \mathbb{C}^{r_2}$ to be the set of those $x = (x_1, \ldots, x_{r_1+r_2})$ in $\mathbb{R}^{r_1} \times \mathbb{C}^{r_2}$ that satisfy:

(i) $|x_{r_1+1} - \bar{x}_{r_1+1}| = 2|\text{Im}(x_{r_1+1})| \leq C\sqrt{|d_K|}$

(ii) $|x_{r_1+1} + \bar{x}_{r_1+1}| = 2|\text{Re}(x_{r_1+1})| \leq \frac{1}{2}$

(iii) $|x_i| < \frac{1}{2}$ if $i \neq r_1 + 1$,

where $C = C(n)$ is chosen to satisfy $\mu(X) > 2^n \mu(D_{Z_K}) = 2^{n-r_2}\sqrt{|d_K|}$. The set $X$ is measurable, symmetric, and convex, so there exists a nonzero $\alpha \in Z_K \cap X$ by the Minkowski Lattice Theorem. Since $|N_{K/Q}(\alpha)|$ is a positive integer, (i), (ii), (iii) imply that $|\text{Im}(\alpha^{(r_1+1)}\sigma_i)| > 1$. Thus $\alpha^{(r_1+1)} \neq \alpha^{(r_1+r_2+1)} = \alpha^{(r_1+1)}$ and $|\alpha^{(r_1+1)}| > 1$. It follows by (i), (ii), (iii) that $\alpha^{(r_1+1)} \neq \alpha^{(i)}$ for any $i \neq r_1 + 1$. Let $L/Q$ be the normal closure of $K/Q$ and extend each $\sigma_i : K \to \mathbb{C}$, $i = 1, \ldots, n$ to $L$. Then $\sigma_{r_1+1}(\alpha) \neq \sigma_i(\alpha)$ for any $i \neq r_1 + 1$. Therefore, $\sigma_{r_1+1}(\alpha) =$
\( \sigma_r^{−1} \sigma_i(\alpha) \) if and only if \( i = j \). It follows that \( \alpha \) has \( n \) distinct conjugates. Since the degree of the minimal polynomial \( m_\mathbb{Q}(\alpha) \) satisfies \( \deg m_\mathbb{Q}(\alpha) = n = [K : \mathbb{Q}] \), we must have \( K = \mathbb{Q}(\alpha) \).

We have, therefore, shown the following:

**Conclusion.** \( K \) is determined by \( m_\mathbb{Q}(\alpha) \) in \( \mathbb{Z}[t] \) (so monic) whose coefficients are elementary symmetric functions in the \( \alpha^{(i)} \).

As \( \alpha \in X \), the coefficients in the conclusion are integers and bounded by a function only depending on \( n \) and \(|d_K|\).

For any fixed \( n_K, d_K \), there exist only finitely many such polynomials, hence only finitely many \( K \) with fixed \( n_K, d_K \) for \( K \) having a complex embedding. By Corollary 88.13, if \( d = d_K \) is fixed, then \( n = n_K \) is bounded. So we are done in this case.

**Case 2.** \( r_2 = 0 \):

Define \( X \subset \mathbb{R}^n \) by those \( x \in \mathbb{R}^n \) satisfying

\[
(i) \quad |x_1| \leq C \sqrt{|d_K|}, \\
(ii) \quad |x_i| \leq \frac{1}{2} \text{ for } j > i
\]

where \( C = C(n) \) is defined as before. Now argue as in Case 1. (This is easier.)

We end by giving a few computations of class numbers using Corollary 88.5.

**Examples 88.15.** 1. Write \( \mathfrak{A} \sim \mathfrak{B} \) if \( \mathfrak{A} \) and \( \mathfrak{B} \) in \( I_{\mathbb{Z}_K} \) are in the same ideal class in \( Cl_{\mathbb{Z}_K} \). Computation involves not only the Minkowski bound, but also knowledge of the way primes split. As example, using the computations in Section 83, we illustrate this for \( K = \mathbb{Q}(\sqrt{-47}) \). We have \( \{1, \frac{1+\sqrt{-47}}{2}\} \) is an integral basis for \( \mathbb{Z}_K \), \( d_K = -47 \), \( B_K = 2\sqrt{47}/\pi = 14/2 < 5 \) and \( N_{K/\mathbb{Q}}(x + y\frac{1+\sqrt{-47}}{2}) = x^2 + xy + 12y^2 \). We know that it suffices to look at prime ideals in \( \mathbb{Z}_K \) lying over 2, 3 respectively. As 2 and 3 split in \( \mathbb{Z}_K \), we have prime factorizations into primes, \( 2\mathbb{Z}_K = \mathfrak{p}_2\mathfrak{p}_2 \) and \( 3\mathbb{Z}_K = \mathfrak{p}_3\mathfrak{p}_3 \) using the computations in Section 83. As 2 is not a norm from \( \mathbb{Z}_K \), we know that \( \mathfrak{p}_2 \not\sim 1 \). We also know that \( 2 \nmid \frac{1+\sqrt{-47}}{2} \) in \( \mathbb{Z}_K \). Changing notation if necessary, this means that \( v_{\mathfrak{p}_2}(\frac{1+\sqrt{-47}}{2}) = 2 \) and \( v_{\mathfrak{p}_2}(\frac{1+\sqrt{-47}}{2}) = 0 \), \( v_{\mathfrak{p}_3}(\frac{1+\sqrt{-47}}{2}) = 1 \) and \( v_{\mathfrak{p}_3}(\frac{1+\sqrt{-47}}{2}) = 0 \), hence \( (\frac{1+\sqrt{-47}}{2})\mathbb{Z}_K \sim \mathfrak{p}_2\mathfrak{p}_2^2 \). It follows that \( \mathfrak{p}_3 \sim \mathfrak{p}_2^2 \) and the class of \( \mathfrak{p}_2 \) generates \( Cl_{\mathbb{Z}_K} \). Since \( \mathfrak{p}_2 \) is not principal, \( h_K > 1 \). Let \( x = 4 + \frac{1+\sqrt{-47}}{2} \), then \( 2 \nmid x \) in \( \mathbb{Z}_K \) and \( N_{K/\mathbb{Q}}(x) = 32 \). Arguing as above, we see that \( \mathfrak{p}_3^5 \) is principal and \( Cl_{\mathbb{Z}_K} \) is cyclic of order 5.
2. We give a real quadratic example: \( K = \mathbb{Q}(\sqrt{31}) \), and again use Section 83. [A fundamental unit is \( \varepsilon = 1520 + 273\sqrt{31} \).] We have \( \mathbb{Z}_K \) has integral basis \( \{1, \sqrt{31}\} \), \( d_K = 4 \cdot 31 \), \( B_K = \sqrt{31} < 6 \), and \( N_{K/\mathbb{Q}}(x + y\sqrt{31}) = x^2 - 31y^2 \). So we need only look at primes lying above 2, 3, 5. We know that 3 and 5 split in \( \mathbb{Z}_K \), say \( 3\mathbb{Z}_K = p_3p_3' \) and \( 5\mathbb{Z}_K = p_5p_5' \). We also know that \( 2\mathbb{Z}_K = p_2^2 \) ramifies. Since \( N_{K/\mathbb{Q}}(39+7\sqrt{31}) = 2 \), we see that \( p_2 \sim 1 \). We also have \( 5 = 6^2 - 31 \cdot 1^2 \) is a norm so \( p_5 \sim 1 \). Moreover, \( -2 \cdot 3 \cdot 5 = -30 = 1^2 - 31 \cdot 1^2 \) is a norm. It follows that \( p_2p_3p_5 \sim 1 \). Therefore, \( p_3 \sim 1 \) also. Hence \( h_K = 1 \).

3. We give the following example that Lang says was Artin’s favorite. Let \( g = t^5 - t + 1 \in \mathbb{Z}[t] \), \( \alpha \) a root of \( g \) in \( \mathbb{C} \) and \( K = \mathbb{Q}(\alpha) \). Let \( L/\mathbb{Q} \) be the normal closure of \( K/\mathbb{Q} \). Then \( g \) is irreducible (as it is irreducible in \( \mathbb{Z}/5\mathbb{Z} \)), and one checks that \( G(L/\mathbb{Q}) \cong S_5 \), \( r_1 = 3 \), \( r_2 = 1 \), and \( \Delta(\alpha) = 2869 = 19 \cdot 151 \) [as \( \Delta(t^5 + at + b) = 5^5b^4 + 2^5a^5 \)]. So \( d_K = \Delta(\alpha) \) (and has only two primes and is square-free) and \( \mathbb{Z}_K = \mathbb{Z}[\alpha] \). We have

\[
B_K \sqrt{|d_K|} = \left( \frac{4}{\pi} \right) \left( \frac{5!}{5^5} \right) \sqrt{2869} < 4.
\]

so every ideal class contains an integral ideal \( \mathfrak{a} \) with \( \mathfrak{N}(\mathfrak{a}) = 1, 2, \) or 3. Suppose that \( \mathfrak{N}(\mathfrak{a}) = 2, 3 \), respectively. Then \( \mathfrak{a} \) is a prime ideal \( \mathfrak{p} \) with \( \mathfrak{N}(\mathfrak{a}) = \mathfrak{p} := 2, 3 \), respectively. It follows that \( f(\mathfrak{p}/\mathfrak{p}\mathbb{Z}) = 1 \) and \( g \mod \mathfrak{p} \) has a root in \( \mathbb{Z}/\mathfrak{p}\mathbb{Z} \), which it does not. Therefore, \( \mathfrak{N}(\mathfrak{a}) = 1 \) and \( \mathbb{Z}_K \) is a PID.

**Exercises 88.16.**

1. Prove Example 88.9.

2. Let \( m \) is a square-free negative number and suppose that \( h_{\mathbb{Z}_Q(\sqrt{-m})} = 1 \). Show all of the following:
   
   (i) We must have \( m \equiv 5 \mod 8 \) except when \( m = -1, -2, \) or \( -7 \).
   
   (ii) If \( p \) is an odd prime satisfying \( m < -4p \), then \( m \) is a nonsquare modulo \( p \).
   
   (iii) If \( m < -19 \), then \( m \) is congruent to one of the following modulo 840: -43, -67, -163, -403, -547, -667.
   
   (iv) If \( -2000 < m < 0 \), then \( \mathbb{Z}_Q(\sqrt{-m}) \) is one of the quadratic number fields given by Example 88.9.

3. Show that \( h_{\mathbb{Z}_Q(\sqrt{-3})} = 2 \).

4. Show that \( h_{\mathbb{Z}_Q(\sqrt{-13})} = 3 \).

5. Show that \( Cl_{\mathbb{Z}_Q(\sqrt{-13})} \) is the Klein 4-group.

6. Show that \( h_{\mathbb{Z}_Q(\sqrt{-2}, \sqrt{-3})} = 1 \). [An integral basis for \( \mathbb{Z}_Q(\sqrt{2}, \sqrt{-3}) \) is \( \{1, \sqrt{2}, (1 + \sqrt{-3})/2, (1 + \sqrt{2} + \sqrt{-6})/2\} \).]
7. Fill in the details to Example 88.15(3).

8. Let $K$ be a number field and $\mathfrak{A} \subset \mathbb{Z}_K$ be an ideal such that $(\mathfrak{A})^m = a\mathbb{Z}_K$ for some $a \in \mathbb{Z}_K$. Show that $\mathfrak{A}\mathbb{Z}_K^{(m\sqrt{a})}$ is principal.

9. Let $K$ be a number field. Show there exists a finite extension $L/K$ such that every nonzero ideal in $\mathbb{Z}_K$ becomes principal in $\mathbb{Z}_L$. 
CHAPTER XVII

Introduction to Commutative Algebra

In this chapter, all rings will be commutative. We shall be particularly interested in developing an algebraic dimension theory that coincides with the intuitive notion of the number of variables needed to describe a geometric object. In addition, we shall establish a generalization of unique factorization into prime ideals as happens in Dedekind domains to arbitrary Noetherian rings.

89. Zariski Topology

Throughout this section, $R$ will denote commutative ring.

For a nonzero commutative ring $R$, we introduced in §40, the Spectrum of $R$,
\[ \text{Spec}(R) := \{ p \mid p < R \text{ a prime ideal} \}, \]
and defined the Zariski topology on it as follows:

If $T$ is a subset of $R$, define
\[ V_R(T) := \{ p \mid p \in \text{Spec}(R) \text{ with } T \subset p \} \]
called the (abstract) variety of $T$. We have

**Lemma 89.1.** Let $R$ be a commutative ring. Then

1. If $T$ is a subset of $R$, then $V_R(T) = V_R(\langle T \rangle)$.
2. If $T_1 \subset T_2$ are subsets of $R$, then $V_R(T_1) \supset V_R(T_2)$.
3. $V_R(\emptyset) = \text{Spec}(R)$.
4. $V_R(R) = \emptyset$.
5. If $T_i, i \in I$, are subsets of $R$, then $V_R(\bigcup_i T_i) = \bigcap_i V_R(T_i)$.
6. If $\mathfrak{A}$ and $\mathfrak{B}$ are ideals in $R$, then
\[ V_R(\mathfrak{A}\mathfrak{B}) = V_R(\mathfrak{A} \cap \mathfrak{B}) = V_R(\mathfrak{A}) \cup V_R(\mathfrak{B}). \]

In particular, the collection
\[ C := \{ V_R(T) \mid T \subset R \} \]
forms a system of closed sets for the Zariski topology on $\text{Spec}(R)$. 599
Proof. This essentially was Exercise 40.15(6).

As an example, we show (6).

Since \( \mathfrak{A} \mathfrak{B} \subset \mathfrak{A} \cap \mathfrak{B} \), we have \( V(\mathfrak{A} \mathfrak{B}) \supset V(\mathfrak{A} \cap \mathfrak{B}) \) by (2).

If \( p \in V(\mathfrak{A} \mathfrak{B}) \), then \( \mathfrak{A} \mathfrak{B} \subset p \), hence \( \mathfrak{A} \subset p \) or \( \mathfrak{B} \subset p \), equivalently, \( p \in V(\mathfrak{A}) \) or \( p \in V(\mathfrak{B}) \), i.e., \( V(\mathfrak{A} \mathfrak{B}) \supset V(\mathfrak{A}) \cup V(\mathfrak{B}) \).

Finally, if \( p \in V(\mathfrak{A}) \cup V(\mathfrak{B}) \), then \( \mathfrak{A} \cap \mathfrak{B} \subset p \), so \( V(\mathfrak{A}) \cup V(\mathfrak{B}) \subset V(\mathfrak{A} \cap \mathfrak{B}) \). □

We shall always assume that \( \text{Spec}(R) \) is given the Zariski topology.

If \( R \) is clear, we shall abbreviate \( V_R(T) \) by \( V(T) \) and if \( \mathfrak{A} = (a_1, \ldots, a_n) \) is a finitely generated ideal in \( R \), we shall write \( V(\mathfrak{A}) \) for \( V(A) \).

The Zariski topology is very coarse. In general, as we shall see, it is not a Hausdorff space, i.e., we cannot necessarily find disjoint open neighborhoods of distinct points. Moreover, we shall see that the set of closed points in \( \text{Spec}(R) \), i.e., those points \( p \) in \( \text{Spec}(R) \) such that \( \{p\} \) is a closed set, is precisely the set of maximal ideals, so points are usually not closed.

Lemma 89.2. Let \( V \) be a subset of \( \text{Spec}(R) \) and \( \overline{V} \) its closure in \( \text{Spec}(R) \). If \( \mathfrak{A} = \bigcap_{p \in V} p \), then \( \overline{V} = V(\mathfrak{A}) \).

Proof. (\( \subset \)) As \( \mathfrak{A} = \bigcap_{V} p \subset \mathfrak{P} \) for all \( \mathfrak{P} \in V \), we have \( V \subset V(\mathfrak{A}) \), hence \( \overline{V} \subset V(\mathfrak{A}) \).

(\( \supset \)) Let \( \overline{V} = V(\mathfrak{B}) \), with \( \mathfrak{B} \) an ideal in \( R \). If \( p \in V \subset \overline{V} = V(\mathfrak{B}) \), then \( \mathfrak{B} \subset p \), hence \( \mathfrak{B} \subset \bigcap_{V} p = \mathfrak{A} \). Therefore, \( \overline{V} = V(\mathfrak{B}) \supset V(\mathfrak{A}) \) by Lemma 89.1(2). □

Corollary 89.3. Let \( R \) be a commutative ring. Then

(1) \( \text{Max}(R) \subset \text{Spec}(R) \) is the set of closed points in \( \text{Spec}(R) \).

(2) Let \( p \) and \( \mathfrak{P} \) be prime ideals. Then \( \mathfrak{P} \in \overline{\{p\}} \), the closure of \( \{p\} \) in \( \text{Spec}(R) \), if and only if \( p \subset \mathfrak{P} \).

(3) \( \overline{\{p\}} = V(p) \) for all prime ideals in \( R \).

(4) \( \text{Spec}(R) \) is a \( T_0 \)-topological space, i.e., if \( p_1 \) and \( p_2 \) are two distinct prime ideals in \( R \), then there exists an open set in \( \text{Spec}(R) \) containing one of \( p_1 \) and \( p_2 \), but not the other.

Proof. (3) follows from the lemma and (2) follows from (3).

(1): Let \( p \in \text{Spec}(R) \). Then we have \( p \) is a closed point if and only if \( \overline{p} = \{q \in \text{Spec}(R) \mid p \subset q\} = \{p\} \) if and only if \( p \in \text{Max}(R) \).

(4): By (2), if \( p_1 \not\subset p_2 \) are prime ideals, then \( p_2 \) lies in the open set \( \text{Spec}(R) \setminus V(p_1) \). □
Definition 89.4. An ideal \( \mathfrak{A} \subset R \) is called a radical ideal if \( \mathfrak{A} = R \) or \( \mathfrak{A} < R \) and
\[
\mathfrak{A} = \sqrt{\mathfrak{A}} := \{ x \in R \mid x^n \in \mathfrak{A} \text{ some } n \in \mathbb{Z}^+ \} = \bigcap_{V(\mathfrak{A})} p = \bigcap_{p \in \text{Spec}(R)} p.
\]
Equivalently, if \( a \) in \( R \) satisfies \( a^n \in \mathfrak{A} \), then \( a \in \mathfrak{A} \).

Examples 89.5. The following are radical ideals:
1. Prime ideals.
2. The nilradical, \( \text{nil}(R) \), of \( R \) as \( \text{nil}(R) = \bigcap_{V(0)} p = \sqrt{(0)} \).

Corollary 89.6. Let \( \mathfrak{A} \) and \( \mathfrak{B} \) be ideals in \( R \). Then \( V(\mathfrak{A}) = V(\sqrt{\mathfrak{A}}) \) and
\[
V(\mathfrak{A}) \subset V(\mathfrak{B}) \text{ if and only if } \sqrt{\mathfrak{A}} \supset \sqrt{\mathfrak{B}}.
\]

Let
\[
\mathcal{R} := \{ \mathfrak{A} \subset R \mid \mathfrak{A} \text{ is a radical ideal} \}
\]
\[
\mathcal{V} := \{ V(\mathfrak{A}) \mid \mathfrak{A} \text{ is an ideal of } R \}.
\]
Then \( V : \mathcal{R} \to \mathcal{V} \) given by \( \mathfrak{A} \mapsto V(\mathfrak{A}) \) is an order-reversing bijection with inverse \( I : \mathcal{V} \to \mathcal{R} \) given by \( V \mapsto \bigcap_{V(A)} p \).

**Proof.** Clearly, \( V(\mathfrak{A}) = V(\sqrt{\mathfrak{A}}) \) and \( V(\mathfrak{A}) \subset V(\mathfrak{B}) \) implies \( \sqrt{\mathfrak{B}} \subset \bigcap_{V(\mathfrak{A})} p = \sqrt{\mathfrak{A}} \).  \( \square \)

We leave the following as an exercise:

**Proposition 89.7.** Let \( r \) be an element of \( R \) and
\[
D(r) := \text{Spec}(R) \setminus V(r) = \{ p \in \text{Spec}(R) \mid r \notin p \}.
\]

Then
1. A finite intersection of basic open sets is a basic open set.
2. \( \{ D(r) \mid r \in R \} \) is a base for the topology of \( \text{Spec}(R) \), i.e., every open set in \( \text{Spec}(R) \) is a union of such \( D(r) \).
3. \( \text{Spec}(R) \) is quasi-compact, i.e., every open cover of \( \text{Spec}(R) \) has a finite subcover.
4. \( \text{Spec}(R) \) is a Hausdorff space if and only if \( \text{Spec}(R) = \text{Max}(R) \).

The set \( D(r) \) above is called a basic open set in \( \text{Spec}(R) \).

**Definition 89.8.** A nonempty topological space is called irreducible if whenever \( U, V \) are nonempty open subsets of \( X \), then \( U \cap V \neq \emptyset \).

A subspace \( Y \) of a topological space \( X \) is called irreducible if it is irreducible in the induced topology.

**Examples 89.9.** 1. Points are irreducible.
2. A nonempty Hausdorff space is irreducible if and only if it consists of a single point.

Proposition 89.10. Let $X$ be a nonempty topological space. Then the following are equivalent:

1. $X$ is irreducible.
2. If $W_1, W_2 < X$ are closed, then $W_1 \cup W_2 < X$.
3. If $\{W_1, \ldots, W_n\}$ is a finite closed cover of $X$, then $X = W_i$ for some $i$.
4. If $U$ is a nonempty open subset of $X$, then $U$ is dense in $X$.
5. Every open set in $X$ is connected.

We leave the proof as an exercise. For varieties, the concept of irreducibility replaces connectivity.

Notation 89.11. We know if $R$ is a nonzero commutative ring, then there exist minimal elements in $\text{Spec}(R)$ by Zorn’s Lemma that we called minimal prime ideals of $R$ (cf. Remark 27.18). We let $\text{Min}(R) := \{p \in \text{Spec}(R) \mid$ there exists no prime ideal $p_0 < p\}$ denote the nonempty set of minimal primes ideals in $R$.

Lemma 89.12. $\text{Spec}(R)$ is irreducible if and only if $\text{nil}(R)$ is a prime ideal if and only if $|\text{Min}(R)| = 1$. In particular, if $R$ is a domain, then $\text{Spec}(R)$ is irreducible.

Proof. We know that $\text{nil}(R) = \bigcap_{\mathfrak{p} \in \text{Spec}(R)} \mathfrak{p} = \bigcap_{\mathfrak{p} \in \text{Min}(R)} \mathfrak{p} \subset \mathfrak{p}$ for every prime ideal $\mathfrak{p}$ and $V(0) = V(\sqrt{0}) = V(\text{nil}(R)) = \text{Spec}(R)$.

Check 89.13. Let $\mathfrak{A}$ and $\mathfrak{B}$ be ideals in $R$. Then

(i) $\mathfrak{A} \subset \text{nil}(R)$ if and only if $\sqrt{\mathfrak{A}} \subset \text{nil}(R)$.

(ii) $\sqrt{\mathfrak{A}\mathfrak{B}} \supset \sqrt{\mathfrak{A}}\sqrt{\mathfrak{B}}$.

Let $\mathfrak{A}$ and $\mathfrak{B}$ be ideals in $R$. Then, using the check, we have $|\text{Min}(R)| = 1$ if and only if $\text{nil}(R)$ is a prime ideal if and only if $\mathfrak{A}\mathfrak{B} \subset \text{nil}(R)$ implies $\mathfrak{A} \subset \text{nil}(R)$ or $\mathfrak{B} \subset \text{nil}(R)$ if and only if $\sqrt{\mathfrak{A}\mathfrak{B}} \subset \text{nil}(R)$ implies $\sqrt{\mathfrak{A}} \subset \text{nil}(R)$ or $\sqrt{\mathfrak{B}} \subset \text{nil}(R)$ if and only if $V(\mathfrak{A}\mathfrak{B}) = V(\mathfrak{A}) \cup V(\mathfrak{B}) = \text{Spec}(R)$ implies $V(\mathfrak{A}) = \text{Spec}(R)$ or $V(\mathfrak{B}) = \text{Spec}(R)$ if and only if $\text{Spec}(R)$ is irreducible.

Since $\text{Spec}(R)$ is a topological space based on an algebraic structure, via the ring $R$, topological maps of spectra should arise from homomorphisms of the underlying algebraic structures. Indeed, we have:
Definition 89.14. Let \( \varphi : A \to B \) be a ring homomorphism of commutative rings. Define the associated map to be the map arising from \( \varphi \) by

\[ a\varphi : \text{Spec}(B) \to \text{Spec}(A) \text{ via } a\varphi(\mathfrak{P}) = \varphi^{-1}(\mathfrak{P}). \]

That this map makes sense follows by:

Lemma 89.15. Let \( \varphi : A \to B \) be a ring homomorphism. Then \( a\varphi(\mathfrak{P}) \) is a prime ideal of \( A \) for every prime ideal \( \mathfrak{P} \) in \( B \), i.e.,

\[ a\varphi : \text{Spec}(B) \to \text{Spec}(A) \]

is defined. Moreover, if \( T \subset A \) is a subset, then

\[ (a\varphi)^{-1}(V(T)) = V(\varphi(T)). \]

In particular, \( a\varphi : \text{Spec}(B) \to \text{Spec}(A) \) is a continuous map.

Proof. If \( \mathfrak{P} < B \) is a prime ideal, then the ideal \( \varphi^{-1}(\mathfrak{P}) < A \) is an ideal and \( \varphi \) induces a monomorphism \( A/\varphi^{-1}(\mathfrak{P}) \to B/\mathfrak{P} \) with \( B/\mathfrak{P} \) a domain. Consequently, \( A/\varphi^{-1}(\mathfrak{P}) \) is also a domain, i.e., \( \varphi^{-1}(\mathfrak{P}) \) is a prime ideal.

We also have

\[ (a\varphi)^{-1}(V(T)) = \{ \mathfrak{P} \in \text{Spec}(B) \mid \varphi^{-1}(\mathfrak{P}) = a\varphi(\mathfrak{P}) \supset T \} \]

\[ = \{ \mathfrak{P} \in \text{Spec}(B) \mid \mathfrak{P} \supset \varphi(T) \} = V(\varphi(T)). \]

It is easy to check:

Remark 89.16. If \( \varphi : A \to B \) and \( \psi : B \to C \) are ring homomorphisms of commutative rings, then \( a(\psi \circ \varphi) = a\varphi \circ a\psi \).

Definition 89.17. A map \( f : X \to Y \) or topological spaces is called dominant if \( \text{im } f \) is dense in \( Y \).

The relationship between \( \varphi \) and \( a\varphi \) is illustrated by:

Lemma 89.18. Let \( \varphi : A \to B \) be a ring homomorphism.

1. If \( \varphi \) is an epimorphism, then \( a\varphi : \text{Spec}(B) \to \text{Spec}(A) \) is an injective closed map with image \( V(\text{ker } \varphi) \). In particular, \( a\varphi \) induces a homeomorphism \( V(\text{ker } \varphi) \cong \text{Spec}(B) \cong \text{Spec}(A/\text{ker } \varphi) \).
2. If \( \mathfrak{A} \) is an ideal, then \( V(\mathfrak{A}) \) is homeomorphic to \( \text{Spec}(R/\mathfrak{A}) \).
3. If \( \varphi \) is a monomorphism, then \( a\varphi : \text{Spec}(B) \to \text{Spec}(A) \) is a dominant map. [It may not be onto.]

Proof. (1) follows from the Correspondence Principle and (1) implies (2).
(3): Suppose that \( a \in A \) satisfies \( D(a) := \text{Spec}(A) \setminus V(a) \) is non-empty. We must show that \( D(a) \cap \text{im}^a \varphi \) is also non-empty. Suppose that this is false, then \( \text{im}^a \varphi \subseteq V(a) \). So

\[
a \in \bigcap_{\text{Spec}(B)}^a \varphi(\mathfrak{p}) = \bigcap_{\text{Spec}(B)} \varphi^{-1}(\mathfrak{p}),
\]

hence \( \varphi(a) \in \bigcap_{\text{Spec}(B)} \mathfrak{p} = \text{nil}(B) \), i.e., \( \varphi(a) \) is nilpotent in \( B \). In particular, there exists a positive integer \( n \) with \( 0 = \varphi(a)^n = \varphi(a^n) \). As \( \varphi \) is monic, \( a^n = 0 \) in \( A \), so \( a \in \text{nil}(A) \subseteq p \) for every prime ideal \( p \). It follows that \( D(a) = \emptyset \), a contradiction. \( \square \)

**Note:** The inclusion \( i : \mathbb{Z} \rightarrow \mathbb{Z}(p) \), with \( \mathbb{Z}(p) = \{ \frac{a}{b} \mid a, b \in \mathbb{Z}, p \nmid b \} \), the localization of \( \mathbb{Z} \) at the prime ideal \( (p) \), is injective, but \( i \) is not surjective as \( \text{Spec}(\mathbb{Z}(p)) = \{ 0, p\mathbb{Z}(p) \} \) (but does have dense image as shown).

Localization will play a decisive role in this chapter as well as being a primary tool in commutative algebra. It is, therefore, useful to coalesce some of the properties of localization, many of which appeared in \( \S 28 \) (some as exercises) and subsequent sections, together with some new observations, leaving all the proofs as useful exercises.

**Remark 89.19.** Let \( R \) be a nonzero commutative ring and \( S \) a multiplicative set in \( R \). We have seen that we have a canonical ring homomorphism \( \varphi_R : R \rightarrow S^{-1}R \) given by \( r \mapsto r/1 \). This makes \( S^{-1}R \) into a commutative \( R \)-algebra. If \( \mathfrak{A} \) in \( R \) is an ideal, then

\[
S^{-1}\mathfrak{A} := \left\{ \frac{a}{s} \mid a \in A, s \in S \right\} \subset S^{-1}R
\]

is an ideal. Let \( \mathfrak{B} \) be another ideal in \( R \). The following are true:

(i) \( S^{-1}\mathfrak{A} = S^{-1}R \) if and only if \( \mathfrak{A} \cap S = \emptyset \).

(ii) \( \mathfrak{A} \subset \varphi_R^{-1}(S^{-1}\mathfrak{A}) \subset R \) is an ideal.

(iii) If \( \mathfrak{D} \) is an ideal in \( S^{-1}R \), then there exists an ideal \( \mathfrak{C} \) in \( R \) satisfying \( S^{-1}\mathfrak{C} = \mathfrak{D} \). Moreover, \( \mathfrak{D} = S^{-1}(\varphi_R^{-1}(\mathfrak{D})) \).

(iv) If \( \mathfrak{A} \subset \mathfrak{B} \) and \( \mathfrak{A} \cap S = \emptyset \), then \( S^{-1}\mathfrak{A} \subset S^{-1}\mathfrak{B} \).

(v) \( S^{-1}(\mathfrak{A} + \mathfrak{B}) = S^{-1}\mathfrak{A} + S^{-1}\mathfrak{B} \).

(vi) \( S^{-1}(\mathfrak{A} \cap \mathfrak{B}) = S^{-1}\mathfrak{A} \cap S^{-1}\mathfrak{B} \).

(vii) \( S^{-1}(\mathfrak{A}\mathfrak{B}) = S^{-1}\mathfrak{A} S^{-1}\mathfrak{B} \).

(viii) The map

\[
\{ p \in \text{Spec} R \mid p \cap S = \emptyset \} \rightarrow \text{Spec}(S^{-1}R) \text{ given by } p \mapsto S^{-1}p
\]

is a bijection and, in fact, a homeomorphism.
As mentioned before, the most important examples of localization of a commutative ring \( R \) are:

(i) \( R_p \) the localization of \( R \) at the multiplicative set \( R \setminus p \) where \( p \) is a prime ideal of \( R \), called the localization at \( p \).

(ii) \( R_a \), the localization of \( R \) at the multiplicative set \( \{a^n \mid n \geq 0\} \) called the localization at \( a \), i.e., at “the open neighborhood \( D(a) \) of \( a \).

We return to the topology that interests us. We have seen that irreducible topological spaces are applicable to our theory. However, connected subspaces are not useful, so we need a substitute for the usual decomposition of a space into connected components into a decomposition of another type. This will be a decomposition into irreducible components.

**Lemma 89.20.** Let \( X \) be a topological space and \( Y \) a subset of \( X \).

1. \( Y \) is irreducible if and only if \( \overline{Y} \) in \( X \) is irreducible.
2. Every irreducible subspace of \( X \) is contained in some maximal irreducible subspace. Any such is closed.

**Proof.** (1): Let \( U \subset X \) be an open subset, then \( Y \subset \overline{U} \) if and only if \( \overline{Y} \subset \overline{U} \). By definition, the closure of \( U \cap \overline{Y} \) in \( \overline{Y} \) is \( U \cap \overline{Y} \), so \( U \cap \overline{Y} \) is dense in \( Y \) if and only if \( U \cap \overline{Y} \) is dense in \( \overline{Y} \).

(2) follows from Zorn’s Lemma. \( \square \)

A maximal irreducible subspace of a nonempty topological space \( X \) is called an irreducible component of \( X \). If \( X \) is a Hausdorff space, then \( \{x\} \), \( x \in X \), are the irreducible components of \( X \), so this is not a useful concept in that case.

**Definition 89.21.** Let \( X \) be a topological space. Then a point \( x \in X \) is called a generic point of \( X \) if \( X = \{x\} \).

**Corollary 89.22.** \( \{V(p) \mid p \in \text{Min}(R)\} \) is the set of irreducible components of \( \text{Spec}(R) \) and \( p \in \text{Spec}(R) \) is a generic point of \( V(p) \).

**Remark 89.23.** Irreducible components may intersect non-trivially. For example, let \( x, y \) be indeterminants, then \( V(xy) \) in \( \text{Spec}(\mathbb{C}[x, y]) \) has irreducible components \( V(x) \) and \( V(y) \) and they intersect in \( (x, y) \), a maximal ideal in \( \mathbb{C}[x, y] \). Topologically, \( V(xy) \cong \text{Spec}(\mathbb{C}[x, y]/(xy)) \).

Identifying \( \text{Max}(\mathbb{C}[x, y]) \) with \( \mathbb{C}^2 \) by \( (x - a, y - b) \mapsto (x, y) \) using the Hilbert Nullstellensatz, we can view the \( y \)-axis as the closed points in \( V(x) \) and the \( x \)-axis as the closed points in \( V(y) \). The origin corresponds to the intersection of \( V(x) \) and \( V(y) \). The following picture illustrates, this where we have identified \( V(x) \cap \text{Max}(\mathbb{C}[x, y]) \) and the
y-axis, \( V(y) \cap \text{Max}(\mathbb{C}[x,y]) \) and the x-axis) and \((x, y)\) with the origin, using the Hilbert Nullstellensatz.

**Definition 89.24.** A topological space \( X \) is called *Noetherian* if the collection of open sets in \( X \) satisfies the ascending chain condition. Equivalently, the collection of closed sets in \( X \) satisfies the descending chain condition.

**Example 89.25.** If \( R \) is a commutative Noetherian ring, then \( \text{Spec}(R) \) is a Noetherian topological space, but the converse is false. For example, let \( F \) be a field, \( A = F[t_1, \ldots, t_n] \), and \( \mathfrak{a} = (t_1, t_2, \ldots, t_n) \). Then the ring \( B = A/\mathfrak{a} \) is not Noetherian, as \( \text{nil}(B) \) is not nilpotent, but \( \text{Spec}(B) \) is a Noetherian space. In fact, \( |\text{Spec}(B)| = 1 \).

**Proposition 89.26.** Let \( X \) be a Noetherian space. Then there exist finitely many irreducible components of \( X \) and \( X \) is their union.

**Proof.** This proof is a typical proof when we have a Noetherian condition. Let

\[
\mathcal{F}(X) := \{ V \mid V \subset X \text{ closed} \}
\]

\[
\mathcal{C}(X) := \{ C \mid C \subset \mathcal{F}(X), C \text{ is a finite union of closed irreducible sets in } X \}\.
\]

**Claim.** \( \mathcal{C}(X) = \mathcal{F}(X) \):

If this is false, then by the Minimal Principle (equivalent to the descending chain condition), there exists an element \( Y \in \mathcal{F}(X) \setminus \mathcal{C}(X) \) that is minimal. As, \( Y \not\in \mathcal{C}(X) \), it is not irreducible, \( Y = Y_1 \cup Y_2 \) with \( Y_i < Y \) and \( Y_i \in \mathcal{F}(X) \), for \( i = 1, 2 \). By minimality, \( Y_i \in \mathcal{C}(X) \), hence \( Y = Y_1 \cup Y_2 \) in \( \mathcal{C}(X) \). This contradiction establishes the claim.

By the claim, we have \( X = \bigcup_{i=1}^{n} C_i \) for some \( C_i \subset X \) closed irreducible. Clearly, we may assume that \( C_j \) is not a subset of \( C_i \) for \( i \neq j \). If \( C \) is an irreducible component of \( X \), then we have \( C \subset \bigcup_{i=1}^{n} (C \cap C_i) \) is irreducible. It follows that \( C \subset C_i \) for some \( i \), hence \( C = C_i \) for
some $i$ as $C$ is a maximal closed irreducible subset of $X$. The result follows. □

**Corollary 89.27.** Let $R$ be a commutative Noetherian ring. Then $\text{Min}(R)$ is a finite set and $\text{Spec}(R) = \bigcup_{\text{Min}(R)} V(p)$.

Our primary interest in this chapter will be the Krull dimension of a ring that we shall just call the *dimension of a ring*. For the convenience of the reader, we define it again, in an equivalent formulation, that will be of considerable interest for us.

**Definition 89.28.** Let $R$ be a nonzero commutative ring and $\mathfrak{p}$ a prime ideal in $R$. We define the *height* of $\mathfrak{p}$ by

$$ht\mathfrak{p} := \sup\{ n \mid p_0 < \cdots < p_n = \mathfrak{p}, \text{ with } p_i \in \text{Spec}(R), i = 1, \ldots, n \}$$

if finite and $= \infty$ if not, i.e., the maximal number of (proper) links of chains of prime ideal in $\mathfrak{p}$. If $\mathfrak{a} < R$ is an ideal, we define the *height* of $\mathfrak{a}$ to be

$$ht\mathfrak{a} := \inf\{ ht p \mid p \in V(\mathfrak{a}) \} = \inf\{ ht p \mid p \in V(\mathfrak{a}) \text{ minimal} \}.$$ 

Then the *(Krull) dimension* of $R$ is defined to be

$$\dim R := \sup\{ ht p \mid p \in \text{Spec}(R) \} = \sup\{ ht m \mid m \in \text{Max}(R) \}$$

if finite and $= \infty$ otherwise.

If $\mathfrak{a}$ is an ideal, we define the *dimension* of $V(\mathfrak{a})$ to be

$$\dim V(\mathfrak{a}) := \dim R/\mathfrak{a}.$$ 

More generally, if $X$ is a topological space, the *combinatorial dimension* of $X$ is defined to be

$$\dim X := \sup\{ n \mid Y_0 < \cdots < Y_n \text{ each } Y_i \subset X \text{ closed and irreducible} \}$$

if finite and $= \infty$ if not. In particular, if $\mathfrak{a} \subset R$ is an ideal, then

$$\dim R = \dim \text{Spec}(R) \text{ and } \dim V(\mathfrak{a}) = \dim \text{Spec}(R/\mathfrak{a}).$$

Note that this is consistent, since $V(\mathfrak{a})$ is homeomorphic to $\text{Spec}(R/\mathfrak{a})$.

**Remarks 89.29.** Let $R$ be a nonzero commutative ring and $F$ a field.

1. If $R$ is a domain, then $\dim R = 0$ if and only if $R$ is a field.
2. $\dim R = 0$ if and only if $\text{Spec}(R) = \text{Max}(R)$.
3. Let $R$ be a PID but not a field. Then $\text{Spec}(R) = \{(0)\} \cup \text{Max}(R)$ with $(0)$ not maximal, so $\dim R = 1$, e.g., $R = \mathbb{Z}$ or $F[t]$. More generally, any Dedekind domain is of dimension one.
4. \( \dim F[t_1, \ldots, t_n] \geq n \), since
\[
0 < (t_1) < (t_1, t_2) < \cdots < (t_1, \ldots, t_n)
\]
is a (proper) chain of prime ideals. In fact, \( \dim F[t_1, \ldots, t_n] = n \), a result that we shall show later.

5. \( \dim F[t_1, \ldots, t_n, \ldots] = \infty \), but there do exist Noetherian rings of infinite dimension.

6. Recall a local ring is a commutative ring with a unique maximal ideal. [Such a ring can still have \( \text{Spec}(R) \) infinite, and in fact, this is the case if \( \dim R > 1 \) and \( R \) is Noetherian.] If \( R \) is a local ring with maximal ideal \( m \), we say \( (R, m) \) is a local ring. We shall see that a local Noetherian ring has finite dimension.

7. Maximal ideals need not have the same height. For example, let \( p = (3) \) in \( \text{Spec}(Z) \) and \( Z_p \) the localization of \( Z \) at \( p \). Then \( Z_p[t] \) is a Noetherian UFD of dimension two containing maximal ideals \( (3, t) \) and \( (3t - 1) \), maximal ideals of height one and two, respectively.

8. Suppose that \( \mathfrak{A} < R \) is an ideal. Then we have:
\[
\text{ht } \mathfrak{A} + \dim R/\mathfrak{A} \leq \dim R,
\]
\[\text{i.e.,}\]
\[
\text{ht } \mathfrak{A} + \dim \text{V}(\mathfrak{A}) \leq \dim \text{Spec}(R).
\]
In general, inequality is possible. If we have equality, then we should view \( \mathfrak{A} \) as \( \text{codim } V(\mathfrak{A}) = \text{codim}_{\text{Spec}(R)} V(\mathfrak{A}) \). One of our primary goals of this chapter is to establish this in the special case when \( R \) is an affine \( F \)-algebra, i.e., a finitely generated commutative \( F \)-algebra, and \( \mathfrak{A} \) is a prime ideal.

9. If \( p \) is a prime ideal in \( R \), let \( p R_p = p_p \), the localization of \( p \) at \( R \setminus p \).
Then \((R_p, p_p)\) is a local ring and satisfies \( \dim R_p = \text{ht } p R_p \).

Remark 89.30. We can also generalize the notions above to modules. Let \( M \) and \( N \) be \( R \)-modules. Then
\[
( N : M ) := \{ x \in R \mid x M \subset N \}
\]
is a ideal in \( R \). E.g., \((0 : M) = \text{ann}_R M \). We define the dimension of \( M \) by
\[
\dim M := \dim V((0 : M)) = \dim V(\text{ann}_R(M))
\]
and the support of \( M \) by
\[
\text{Supp}(M) := \{ p \in \text{Spec}(R) \mid M_p \neq 0 \},
\]
where \( M_p = \{ m/s \mid m \in M, s \in R \setminus p \} \), the localization of \( M \) at \( p \), an \( R_p \)-module. [We leave the proof that we can localize an \( R \)-module to an \( S^{-1}R \)-module with \( S \) a multiplicative set in \( R \) as an exercise.]
Exercises 89.31.
1. Prove Proposition 89.7.
2. Prove Proposition 89.10.
3. Show Check 89.13 is true.
4. Check Remark 89.16.
5. Prove the assertions in Remark 89.19.
6. Let $R$ be a nonzero commutative ring, $S$ a multiplicative set in $R$, and $M$ an $R$-module. Show that $S^{-1}M := \{\frac{r}{s} \mid r \in R, s \in S\}$ with the obvious definition is an $S^{-1}R$-module called the localization of $M$ at $S$. If

$$0 \to M' \xrightarrow{f} M \xrightarrow{g} M'' \to 0$$

is an exact sequence of $R$-modules, prove that the sequence

$$0 \to S^{-1}M' \xrightarrow{S^{-1}f} S^{-1}M \xrightarrow{S^{-1}g} S^{-1}M'' \to 0$$

of $S^{-1}R$-modules is exact where, e.g., $S^{-1}f(\frac{r}{s}m) = \frac{r}{s}f(m)$, i.e., the $S^{-1}R$-module homomorphism induced by $f$.

7. In the previous exercise, if $S = R \setminus \mathfrak{p}$ with $\mathfrak{p}$ a prime ideal in $R$, and $f : M \to N$ is an $R$-homomorphism, let $f_\mathfrak{p}$ denote $S^{-1}f$. Then show ker$(f_\mathfrak{p}) = (\text{ker } f)_\mathfrak{p}$ and coker$(f_\mathfrak{p}) = (\text{coker } f)_\mathfrak{p}$ for all prime ideals $\mathfrak{p}$ in $R$. Moreover, show that the following are equivalent:

(i) $0 \to M' \xrightarrow{f} M \xrightarrow{g} M'' \to 0$ is an exact sequence of $R$-modules.

(ii) $0 \to M'_\mathfrak{p} \xrightarrow{f_\mathfrak{p}} M_\mathfrak{p} \xrightarrow{g_\mathfrak{p}} M''_\mathfrak{p} \to 0$ is an exact sequence of $R_\mathfrak{p}$-modules for all prime ideals $\mathfrak{p}$.

(iii) $0 \to M'_m \xrightarrow{f_m} M_m \xrightarrow{g_m} M''_m \to 0$ is an exact sequence of $R_m$-modules for all maximal ideals $m$.

In particular, $f$ is an $R$-monomorphism (respectively, $R$-epimorphism) if and only if $f_\mathfrak{p}$ is an $R_\mathfrak{p}$-monomorphism (respectively $R_\mathfrak{p}$-epimorphism) for all prime ideals $\mathfrak{p}$ in $R$ if and only if $f_\mathfrak{p}$ is an $R_\mathfrak{m}$-monomorphism (respectively $R_\mathfrak{m}$-epimorphism) for all maximal ideals $\mathfrak{m}$ in $R$. It follows that $M = 0$ if and only if $M_\mathfrak{p} = 0$ for all prime ideals $\mathfrak{p}$ in $R$ if and only if $M_m = 0$ for all maximal ideals $\mathfrak{m}$ in $R$.

8. Let $R$ be a nonzero commutative ring and $a_1, \ldots, a_n \in R$. Show $D(a_1 \cdots a_n) = \bigcap_{i=1}^n D(a_i)$ and every open set in Spec$(R)$ is a union of basic open sets.

9. Show the assertions in Example 89.25 are true.

10. Let $X$ be a topological space. Show that $X$ is Noetherian if and only if every open subset of $X$ is quasi-compact.
11. For \( f \in R \), let \( X = \text{Spec}(R) \) and \( X_f = \text{Spec}(R_f) \), where \( R_f \) the localization of \( R \) at the basic open set \( D(f) \). Write \( R(X_f) \) for the localization \( R_f \).

For \( f, g, h \in R \) and \( U = X_f, U' = X_g, U'' = X_h \), show the following:

(a) \( R(U) \) depends only on \( U \) and not on \( f \).

(b) If \( U' \subset U \), then there exists a positive integer \( n \) and an element \( x \in R \) such that \( g^n = xf \). Using this, defines a homomorphism \( \rho_{UU'} : R(U) \to R(U') \) via \( a/f^m \mapsto ax^m/g^m \). This map depends only on \( U \) and \( U' \) and is called the restriction homomorphism.

(c) If \( U = U' \), then \( \rho_{UU} \) is the identity map.

(d) If \( U'' \subset U' \subset U \), then the diagram

\[
\begin{array}{ccc}
R(U) & \xrightarrow{\rho_{UU'}} & R(U') \\
\downarrow{\rho_{UU''}} & & \downarrow{\rho_{U'U''}} \\
R(U'') & & \\
\end{array}
\]

commutes. This implies that \( X := (X, R) \) is a presheaf of rings.

(e) Let \( \{U_i \mid i \in I\} \) be a finite covering of a basic open set \( U \) in \( X \) by basic open sets. Suppose that \( a \in U, a_i \in U_i \) for each \( i \). Then

(i) If \( a \in R(U) \) satisfies \( \rho_{UU_i}(a) = 0 \) for all \( i \), then \( a = 0 \).

(ii) If \( \rho_{U_i(U_i \cap U_j)}(a_i) = \rho_{U_j(U_i \cap U_j)}(a_j) \) in \( R(U_i \cap U_j) \) for all \( i, j \), then there exists an element \( r \in R(U) \) satisfying \( \rho_{UU_i}(r) = a_i \) for all \( i \).

This implies that \( X \) is a sheaf of rings. We call \( (X, R) \) an affine scheme.

[ If \( x = p \), then the local ring \( R_p \) is the stalk of \( X \) at \( p \). It is the set of germs at \( p \) on \( X \), i.e., equivalence classes of elements \( r \) where \( r \) is defined in some \( R(U) \), \( U \) a basic open set, and if \( r' \) is defined in \( R(U') \), then \( r \sim r' \) if there exists a basic open neighborhood \( V \) of \( p \) in \( U \cap U' \) such that \( \rho_{UV}(r) = \rho_{UV}(r') \) in \( R(V) \).]

90. Integral Extensions of Commutative Rings

As is true in this whole chapter, \( R \) is a commutative ring.

In this section, we look further at the notion of integral extensions. We use the material studied in Sections 76 and 77.

**Notation 90.1.** We begin with some new notation. If \( A \) is a commutative \( R \)-algebra via the ring homomorphism \( \varphi_A : R \to A \), we shall
write \( ra \) for \( \varphi(r)a, \) \( r \in R, a \in A. \) If \( a_1, \ldots, a_n \) are elements of \( A, \) then we shall write \( R[a_1, \ldots, a_n] \) for \( \varphi(R)[a_1, \ldots, a_n] \) (even if \( \varphi \) is not a monomorphism). If \( S \subset R \) is a multiplicative set, we shall write \( S^{-1}A \) for \( (\varphi(S))^{-1}A. \) So if \( r \in R, s \in S, \) and \( a \in A, \) we write \( \frac{r}{s}a \) for \( \frac{\varphi(r)}{\varphi(s)}a. \)

If \( p \) is a prime ideal in \( R, \) we write \( A_p \) for \( S^{-1}A \) where \( S = R \setminus p. \)

**Definition 90.2.** Let \( \varphi : R \to A \) be a ring homomorphism of commutative rings, we shall view \( A \) as an \( R \)-algebra via \( \varphi \) unless otherwise stated. We say

1. \( \varphi \) is of **finite type** if \( A \) is a finitely generated \( R \)-algebra.
2. \( \varphi \) is **finite** if \( A \) is a finitely generated \( R \)-module (via \( \varphi \)).
   
   [One also calls a finitely generated \( R \)-module a **finite \( R \)-module**.]
3. \( \varphi \) is **integral** if \( A/\varphi(R) \) is an integral extension.

We leave the various remarks about ring homomorphisms of commutative rings and localization as useful exercises for the reader.

**Remarks 90.3.** Let \( \varphi : R \to A \) be a ring homomorphism of commutative rings.

1. \( \varphi \) is a finite map if and only if it is integral and finite type. (Cf. Exercise 76.15(4).)
2. If \( \varphi \) is an epimorphism, then \( \varphi \) is finite, hence integral.
3. If \( \varphi \) is integral and \( \mathfrak{A} \) an ideal in \( A, \) then the induced map \( \tilde{\varphi} : R/\varphi^{-1}(\mathfrak{A}) \to A/\mathfrak{A} \) is integral and injective.
4. The composition of ring homomorphisms of finite type (respectively, finite, integral) are of finite type (respectively, finite, integral).
5. Suppose that \( \varphi \) is the inclusion and \( \mathfrak{A} < R \) an ideal. As before, we write \( A/R. \) It is possible that \( \mathfrak{A}A = A \) and \( \mathfrak{A} < (\mathfrak{A}A) \cap R. \) For example, \( (2) \) is a prime ideal in \( \mathbb{Z} \) but \( (2)\mathbb{Q} = \mathbb{Q}. \) We shall see that this does not happen if \( A/R \) is integral.
6. Suppose that \( \varphi \) is the inclusion with \( R \) a domain and \( K = A \) a field containing \( R. \) As in §76, we shall let \( R_K \) denote the integral closure of \( R \) in \( K. \) In this chapter, if \( R \) is integrally closed in its quotient field \( F, \) i.e., \( R = R_F, \) we shall use geometric language and say that \( R \) is a **normal domain** instead of saying \( R \) is integrally closed.

**Remarks 90.4.** Let \( \varphi : R \to A \) be a ring homomorphism of commutative rings and \( S \) a multiplicative subset of \( R. \)

1. If \( \varphi \) is a monomorphism, so is \( \varphi_{S^{-1}R} : S^{-1}R \to S^{-1}A. \) In particular, this applies if \( \varphi \) is the inclusion map. We then view \( \varphi_{S^{-1}R} \) as the
inclusion. If, in addition, \( A = K \) is a field \( qf(R) \) and \( \mathfrak{p} \) a prime ideal in \( R \), then \( (R_{\mathfrak{p}})_{K} = (R_{K})_{\mathfrak{p}} \).

2. If \( \varphi \) is of finite type (respectively, finite, integral), then so is \( \varphi_{S^{-1}R} \) in which case the integral closure of \( S^{-1}R \) in \( S^{-1}A \) is \( S^{-1}C \) where \( C \) is the integral closure of \( R \) in \( A \).

In the more elementary part of this tome, we proved the Hilbert Nullstellsatz. We wish to give deeper insight into why this theorem (rather than the proof we gave) is true. One of the keys to the proof, essentially an equivalent formulation, was Zariski’s Lemma 40.10 whose proof we shall derive again later. To do so we must look carefully at integral extensions. We begin with a summary and the methodology that we shall use.

**Summary 90.5.** Let \( i : A \subset B \) be the inclusion of two commutative rings and \( S \subset A \) a multiplicative set.

1. \( S^{-1}i : S^{-1}A \to S^{-1}B \) is a ring monomorphism, that we view as an inclusion.

2. Let \( S^{-1}\varphi_{A} : A \to S^{-1}A \) by \( a \to \frac{a}{1} \), be the canonical homomorphism. Then \( a\varphi_{A} : \text{Spec} S^{-1}A \to \{ \mathfrak{p} \in \text{Spec}(A) \mid \mathfrak{p} \cap S = \emptyset \} \) is a homeomorphism. In particular, if \( \mathfrak{p} \) is a prime ideal in \( A \) such that \( \mathfrak{p} \cap S = \emptyset \), then \( \mathfrak{p} = (S^{-1}\mathfrak{p}) \cap A := a\varphi_{A}(S^{-1}\mathfrak{p}) \).

3. If \( \mathfrak{p} \) is a prime ideal in \( A \), then \( (A_{\mathfrak{p}}, \mathfrak{p}A_{\mathfrak{p}}) \) is a local ring.

4. A major methodology for trying to prove commutative algebra results is by:

   (a) **Localize:** If \( \mathfrak{p} \) is a prime ideal, then \( A_{\mathfrak{p}} \subset B_{\mathfrak{p}} \) with \( (A_{\mathfrak{p}}, \mathfrak{p}A_{\mathfrak{p}}) \) a local ring. Then try to prove the results here and pullback.

   (b) **Quotient:** If \( \mathfrak{P} \) is a prime ideal in \( B \), let \( \mathfrak{p} = \mathfrak{P} \cap A \) a prime ideal in \( A \) with the induced map \( A/\mathfrak{p} \to B/\mathfrak{P} \) a homomorphism of domains. Then try to prove the results here and pullback.

**Lemma 90.6.** Let \( A \) be a commutative \( R \)-algebra and \( u \) a unit in \( A \). Then \( u^{-1} \) is integral over \( R \) if and only if \( u^{-1} \) lies in \( R[u] \).

**Proof.** The element \( u^{-1} \) is integral over \( R \) if and only if there exists an equation

\[
u^{-n} + r_1u^{-n+1} + \cdots + r_n = 0 \quad \text{in} \ A \quad \text{for some} \ r_1, \ldots, r_n \in R
\]

if and only if

\[
u(r_1 + r_2u + \cdots + r_nu^{n-1}) = -1 \quad \text{in} \ A \quad \text{for some} \ r_1, \ldots, r_n \in R
\]

if and only if \( u^{-1} \) lies in \( R[u] \).

**Proposition 90.7.** Let \( A \subset B \) be domains with \( B/A \) integral. Then \( A \) is a field if and only if \( B \) is a field.
Proof. ($\Leftarrow$): Let $x \in A$ be nonzero. Then $x^{-1}$ lies in the field $B$ and is integral over $A$ if and only if $x^{-1} \in A[x] = A$ by the lemma.

($\Rightarrow$): Let $y \in B$ be nonzero. Then there exists $a_1, \ldots, a_n$ in $A$ such that $y^n + a_1y^{n-1} + \cdots + a_n = 0$ in $B$ for some $n$, as $B/A$ is integral. Since $B$ is a domain, we may assume that $a_n$ is nonzero. Therefore, $a_n^{-1}$ lies in $A \subset qf(B)$, the quotient field of $B$, as $A$ is a field. Consequently, $y^{-1} = -a_n^{-1}(y^{n-1} + \cdots + a_{n-1})$ in $qf(B)$, hence lies in $A[y] \subset B$. □

Corollary 90.8. Let $\varphi : A \to B$ be an integral ring homomorphism, $\mathfrak{P}$ a prime ideal in $B$. Then $\mathfrak{P}$ is a maximal ideal in $B$ if and only if $a\varphi(\mathfrak{P})$ is a maximal ideal in $A$. In particular, the restriction of $a\varphi$ to $\text{Max}(B)$ gives the map $a\varphi|_{\text{Max}(B)} : \text{Max}(B) \to \text{Max}(A)$.

Proof. Let $\mathfrak{P}$ be a prime ideal in $B$. Then $\varphi$ induces a monomorphism $\overline{\varphi} : A/\varphi^{-1}(\mathfrak{P}) \to B/\mathfrak{P}$ which is integral by Remark 90.3(3). □

Definition 90.9. Let $R$ be a commutative ring. The Jacobson radical of $R$ is defined to be $\text{rad}(R) := \bigcap_{\mathfrak{m} \in \text{Max}(R)} \mathfrak{m}$.

One of the most useful results is the following lemma, apparently first proved independently by Azumaya and Krull:

Lemma 90.10. (Nakayama’s Lemma) Let $\mathfrak{A}$ be an ideal in $\text{rad}(R)$ and $M$ a finitely generated $R$-module. If $M = \mathfrak{A}M$, then $M = 0$.

Proof. Suppose that $M \neq 0$ and $M = \sum_{i=1}^n Rm_i$ with $n$ minimal. Then there exist $a_1, \ldots, a_n$ in $\mathfrak{A}$ satisfying $m_1 = \sum_{i=1}^n a_im_i$, as $M = \mathfrak{A}M$. Thus $(1 - a_1)m_1 = \sum_{i=2}^n a_im_i$. Since each $a_i$ lies in $\text{rad}(R) \subset \mathfrak{m}$ for every maximal ideal $\mathfrak{m}$ in $R$, we have $1 - a_1 \notin \mathfrak{m}$ for every maximal ideal $\mathfrak{m}$, hence is a unit in $R$. It follows that $m_1$ lies in $\sum_{i=2} Rm_i$, contradicting the minimality of $n$. □

The following consequence is also called Nakayama’s Lemma.

Corollary 90.11. Let $\mathfrak{A} \subset \text{rad}(R)$ be an ideal in $R$, $M$ a finitely generated $R$-module, and $N \subset M$ a submodule. If $M = N + \mathfrak{A}M$, then $M = N$.

Proof. We have $M/N$ is a finitely generated $R$-module satisfying $M/N = \mathfrak{A}(M/N)$, so $M/N = 0$ by Nakayama’s Lemma. Hence $M = N$. □

We leave the proof of the next useful corollary as an exercise.
Corollary 90.12. Let \((R, \mathfrak{m})\) be a local ring, \(M\) a finitely generated \(R\)-module, \(\pi : M \to M/\mathfrak{m}\), the canonical epimorphism,
\[
S = \{x_1, \ldots, x_n\}, \quad \text{and} \quad \overline{S} = \{\overline{x_1}, \ldots, \overline{x_n}\}.
\]
Then
1. \(S\) generates \(M\) if and only if \(\overline{S}\) spans the \(R/\mathfrak{m}\)-vector space \(\overline{M} = M/\mathfrak{m}M\).
2. \(S\) is a minimal generating set (obvious definition) for \(M\) if and only if \(\overline{S}\) is an \(R/\mathfrak{m}\)-basis for \(\overline{M}\) (and \(\dim_{R/\mathfrak{m}} \overline{M} = n\)).

Remark 90.13. We shall see later that if \((R, \mathfrak{m})\) is a Noetherian local ring that
\[
\dim R = \text{ht} \mathfrak{m} \leq \dim_{R/\mathfrak{m}} \mathfrak{m}/\mathfrak{m}^2 < \infty.
\]
If we have equality in the above, then \((R, \mathfrak{m})\) is called a regular local ring. This notion of regular is the algebraic replacement for the geometric concept of non-singularity at a point.

A basic theorem in the study of integral extensions, apparently originally due to Krull, is the next result that we now prove. The proof is an excellent illustration of the methods mentioned above.

Theorem 90.14. (Cohen-Seidenberg Theorem)
Let \(B/A\) be integral. Then:
1. (Incomparability) If \(\mathfrak{p}_1 \subsetneq \mathfrak{p}_2\) are prime ideals in \(B\), and satisfy \(\mathfrak{p}_1 \cap A = \mathfrak{p}_2 \cap A\), then \(\mathfrak{p}_1 = \mathfrak{p}_2\).
2. (Lying Over) \(i : \text{Spec}(B) \to \text{Spec}(A)\) is surjective, where \(i : A \to B\) is the inclusion map, i.e., (in the notation of Section 79 that we shall continue to use) if \(\mathfrak{p}\) is a prime ideal in \(A\), there exists a prime ideal \(\mathfrak{q}\) in \(B\) lying over \(\mathfrak{p}\).
3. (Going Up). If \(\mathfrak{p}_1 \subsetneq \mathfrak{p}_2\) are prime ideals in \(A\) and \(\mathfrak{p}_1\) is a prime ideal in \(B\) lying over \(\mathfrak{p}_1\), then there exists a prime ideal in \(V(\mathfrak{p}_1)\) lying over \(\mathfrak{p}_2\).

\[
\begin{array}{c}
\mathfrak{p}_1 \subset \mathfrak{p}_2 \\
\mathfrak{p}_1 \subset \mathfrak{p}_2
\end{array}
\]

(4) \(\dim A = \dim B\).

**Proof.** (1): Let \(\mathfrak{p} = \mathfrak{p}_1 \cap A = \mathfrak{p}_2 \cap A\) in \(\text{Spec}(A)\) and \(S = A \setminus \mathfrak{p}\). Then \(S \cap \overline{\mathfrak{p}_i} = \emptyset\) for \(i = 1, 2\), so \(S^{-1}B/S^{-1}A\) is integral, \((S^{-1}A, S^{-1}\mathfrak{p})\) is a local ring with
\[
S^{-1}\mathfrak{p}_i \cap S^{-1}A = S^{-1}(\mathfrak{p}_i \cap A) = S^{-1}\mathfrak{p}.
\]
Therefore, $S^{-1}\mathfrak{p}_i \in \text{Spec}(S^{-1}B)$ lies over $S^{-1}\mathfrak{p}$, for $i = 1, 2$. By Corollary 90.8, $S^{-1}\mathfrak{p}_i$ is a maximal ideal in $S^{-1}B$ for $i = 1, 2$. It follows that $\mathfrak{p}_1 \subset \mathfrak{p}_2$ implies $S^{-1}\mathfrak{p}_1 = S^{-1}\mathfrak{p}_2$. Consequently, $\mathfrak{p}_1 = \mathfrak{p}_2$ as $\mathfrak{p}_i \cap S = \emptyset$ for $i = 1, 2$.

(2): Let $S = A \setminus \mathfrak{p}$, then $S^{-1}B/S^{-1}A$ is integral. By the diagram

$$
\begin{array}{ccc}
B & \longrightarrow & S^{-1}B \\
\downarrow & & \downarrow \\
A & \longrightarrow & S^{-1}A
\end{array}
$$

if $S^{-1}\mathfrak{p}$ lies over $S^{-1}\mathfrak{p}$, we see that it suffices to replace $A$ by $S^{-1}A$ and assume that $(A, \mathfrak{p})$ is a local ring with $B/A$ integral. In particular, it suffices to show that there exists a prime ideal $\mathfrak{p}$ in $B$ lying over $\mathfrak{p}$. Suppose that $\mathfrak{p}B < B$. Then there exists an ideal $\mathfrak{m} \in V(\mathfrak{p}B) \cap \text{Max}(B)$. Since $B/A$ is integral, $\mathfrak{m} \cap A \in \text{Max}(A) = \{\mathfrak{p}\}$, i.e., $\mathfrak{m} \cap A = \mathfrak{p}$ and we would be done. So we may assume that $\mathfrak{p}B = B$. We then have an equation

$$
1 = \sum_{i=1}^{m} p_i b_i \text{ for some } p_i \in \mathfrak{p}, \ b_i \in B.
$$

Set $B' = A[b_1, \ldots, b_m]$ a finitely generated commutative $A$-algebra. As $B/A$ is integral, $B'$ is a finitely generated $A$-module satisfying $\mathfrak{p}B' = B'$. It follows by Nakayama’s Lemma, that $B' = 0$, which is impossible. Hence $\mathfrak{p}B < B$ and (2) follows.

(3): Let $\mathfrak{p}_1 \in \text{Spec}(B)$ lie over $\mathfrak{p}_1 \in \text{Spec}(A)$. Since the inclusion $A \subset B$ induces an integral monomorphism $A/\mathfrak{p}_1 \rightarrow B/\mathfrak{p}_1$, by Lying Over and the Correspondence Principle, there exists a prime ideal $\mathfrak{p}_2$ in $V(\mathfrak{p}_1)$ with $\mathfrak{p}_2/\mathfrak{p}_1$ lying over $\mathfrak{p}_2/\mathfrak{p}_1$ in $\text{Spec}(A/\mathfrak{p}_1)$. Therefore, $\mathfrak{p}_2$ lies over $\mathfrak{p}_1$ by the Correspondence Principle and (3) follows.

(2): Let

$$
\mathfrak{p}_1 \subset \cdots \subset \mathfrak{p}_n
$$

be a chain of prime ideals in $B$. By Incomparability,

$$
\mathfrak{p}_1 \cap A \subset \cdots \subset \mathfrak{p}_n \cap A
$$

is a chain of prime ideals in $A$, so $\dim B \leq \dim A$. If

$$
\mathfrak{p}_1 \subset \cdots \subset \mathfrak{p}_n
$$

is a chain of prime ideals in $A$, then by Lying Over and Going Up, we can construct a chain of prime ideals

$$
\mathfrak{p}_1 \subset \cdots \subset \mathfrak{p}_n
$$
in $B$, so $\dim A \leq \dim B$. Therefore, $\dim A = \dim B$. □

In the theorem, we needed the condition that $A \subset B$, i.e., that we have a monomorphism $A \to B$. Indeed a counterexample is provided by the canonical epimorphism $\mathbb{Z} \to \mathbb{Z}/n\mathbb{Z}$, $n > 1$ if not.

**Corollary 90.15.** Let $B/A$ be an integral extension with $B$ a domain and $\mathfrak{B} \subset B$ a nonzero ideal. Then $\mathfrak{B} \cap A < A$ is a nonzero ideal.

**Proof.** Let $\mathfrak{P} \in V(\mathfrak{B})$. Certainly, $\mathfrak{P} \cap A < A$ as $1/\in \mathfrak{P}$. Since $A$ and $B$ are domains, the zero ideal in $B$ lies over the zero ideal in $A$. By Incomparability, $\mathfrak{P}$ does not lie over $(0)$. Suppose that $\mathfrak{B} \cap A = 0$. Then the induced monomorphism $A \to B/\mathfrak{B}$ is integral. Therefore, by Lying Over and the Correspondence Principle, there exists a prime ideal $\mathfrak{P} \in V(\mathfrak{B})$ satisfying $\mathfrak{P}/\mathfrak{B}$ lies over $(0)$ in $\text{Spec}(A)$. As $\mathfrak{B} \cap A = 0$, this means that $\mathfrak{P}$ lies over $(0)$, a contradiction. □

We need to generalize Proposition 80.1, whose former proof can be adapted to this generalization (essentially by localization). Moreover, we shall give a different proof that uses one of the most useful results, called the Prime Avoidance Lemma that we previously left as an exercise, but now prove.

**Lemma 90.16.** (Prime Avoidance Lemma) Let $\mathfrak{A}_1, \ldots, \mathfrak{A}_n$ be ideals in $R$, at least $n - 2$ of which are prime. Let $S \subset R$ be a subrng (it does not have to have a 1) contained in $\mathfrak{A}_1 \cup \cdots \cup \mathfrak{A}_n$. Then there exists a $j$ such that $S \subset \mathfrak{A}_j$. In particular, if $\mathfrak{P}_1, \ldots, \mathfrak{P}_n$ are prime ideals in $R$ and $\mathfrak{B}$ is an ideal properly contained in $S$ satisfying $S \setminus \mathfrak{B} \subset \mathfrak{P}_1 \cup \cdots \cup \mathfrak{P}_n$, then $S$ lies in one of the $\mathfrak{P}_i$'s.

**Proof.** The case $n = 1$ is trivial, so assume that $n > 1$ and assume that the result is false. By induction, $S \not\subset \mathfrak{A}_1 \cup \cdots \mathfrak{A}_j \cup \cdots \cup \mathfrak{A}_n$ for each $i$ where $\widehat{\cup}$ means omit. So for each $i = 1, \ldots, n$, there exists an element $x_i \in S \setminus \mathfrak{A}_1 \cup \cdots \mathfrak{A}_i \cup \cdots \cup \mathfrak{A}_n$. Therefore, we have $x_i \in \mathfrak{A}_i$ for every $i = 1, \ldots, n$ and $x_i \not\in \mathfrak{A}_j$ for all $j \neq i$. If $n > 2$, we may assume that $\mathfrak{A}_1$ is a prime ideal. Let $y = x_1 + x_2 \cdots x_n \in S$. If $n = 2$, then $y = x_1 + x_2$ does not lie in $\mathfrak{A}_1 + \mathfrak{A}_2$, a contradiction. So we may assume that $n > 2$. As $\mathfrak{A}_1$ is a prime ideal and $x_i \not\in \mathfrak{A}_1$ for all $i > 1$, we have $x_2 \cdots x_n \not\in \mathfrak{A}_1$. Consequently, $y \not\in \mathfrak{A}_1 \cup \cdots \cup \mathfrak{A}_n$, a contradiction. □

Our generalization of Proposition 80.1 needs us to replace the condition of a finite Galois extension of fields by an arbitrary normal extension of fields. This requires that we extend the definition of the norm of a finite separable field extension. Let $L/F$ be a finite normal extension of fields, and $K$ the separable closure of $F$ in $K$. We call
Define the norm $N_{L/F} : L \to F$ by $N_{L/F}(x) = (\prod_{G(L/F)} \sigma(x))^{[L:F]}$, i.e., $N_{L/K}(x) = N_{K/F}(x^{[L:F]})$. Then the norm still satisfies all the expected properties.

**Theorem 90.17.** Let $A$ be a normal domain (i.e., an integrally closed domain) with quotient field $F$ and $L/F$ a normal extension of fields. Let $\mathfrak{p}_1$ and $\mathfrak{p}_2$ be prime ideals in $A_L$ lying over $p$ in Spec$(A)$, i.e., $\mathfrak{p}_1 \cap A = p = \mathfrak{p}_2 \cap A$. Then there exists an element $\sigma$ in $G(L/F)$, the Galois group of $L/F$, satisfying $\mathfrak{p}_1 = \sigma(\mathfrak{p}_2)$, i.e., if $i : A \to A_L$ is the inclusion map, then $G(L/F)$ acts transitively on the fibers of $^a i : \text{Spec}(A_L) \to \text{Spec}(A)$.

**Proof.** If $b \in A_L$ and $\sigma \in G(L/F)$, then $\sigma(b) \in L$ is integral over $\sigma(A) = A$, i.e., $\sigma(b) \in (\sigma(A))_L = A_L$, so $\sigma(A_L) = A_L$. As $\sigma^{-1}$ lies in $G(L/F)$, the map $\sigma : A_L \to A_L$ is an $A$-algebra isomorphism. If $\mathfrak{p}$ is a prime ideal in $A_L$ and $\sigma \in G(L/F)$, then $\sigma(\mathfrak{p}) \in \text{Spec}(A_L)$ and $\mathfrak{p} \cap A = \sigma(\mathfrak{p} \cap A) = \sigma(\mathfrak{p}) \cap A$. So every prime ideal in the orbit $G(L/F)\mathfrak{p} := \{\sigma(\mathfrak{p}) \mid \sigma \in G(L/F)\}$ lies in the fiber $^{(a)^{-1}}(p)$ of $p$.

**Case 1.** $L/F$ is finite.

Suppose $\mathfrak{p}_2$ lies over $\mathfrak{p}_1 \cap A$. To show that $\mathfrak{p}_2 \in G(L/F)\mathfrak{p}_1$. Suppose not. By Incomparability, $\mathfrak{p}_2 \not\subseteq \sigma(\mathfrak{p}_1)$ for any $\sigma \in G(L/F)$. By the Prime Avoidance Lemma, as $G(L/F)\mathfrak{p}_1$ is finite, $\mathfrak{p}_2 \not\subseteq \bigcup_{\sigma \in G(L/F)} \sigma(\mathfrak{p}_1)$. Let $a$ be an element in $\mathfrak{p}_2 \setminus \bigcup_{\sigma \in G(L/F)} \sigma(\mathfrak{p}_1)$. Then $\sigma(a) \notin \mathfrak{p}_1$ for any $\sigma \in G(L/F)$. As $A$ is a normal domain, we have

$$N_{L/F}(a) = \left( \prod_{\sigma \in G(L/F)} \sigma(a) \right)^{[L:F]} \text{ lies in } A_L \cap F = A_F = A.$$ 

Therefore,

$$N_{L/F}(a) \in \mathfrak{p}_2 \cap A = p = \mathfrak{p}_1 \cap A \subset \mathfrak{p}_1.$$ 

As $\mathfrak{p}_1$ is a prime ideal and $G(L/F)$ is a finite group, there exists an $\sigma$ in $G(L/F)$ satisfying $\sigma(a) \in \mathfrak{p}_1$, a contradiction.

**Case 2.** $L/F$ is possibly infinite.

This is left as an exercise (use Zorn’s Lemma, since $L$ is a union of finite normal extensions).

**Corollary 90.18.** Let $A$ be a normal domain with quotient field $F$ and $K/F$ a finite field extension. Then $^a i : \text{Spec}(A_K) \to \text{Spec}(A)$ has finite fibers, where $i$ is the inclusion map of $A$ in $A_K$, i.e., $^{(a)^{-1}}(p)$ is a finite set for all prime ideals $p$ in $A$. 


A ring homomorphism whose associated map has finite fibers is called *quasi-finite*.

**Proof.** Let $L/F$ be a normal closure of $K/F$. Then $L/F$ is a finite normal extension, hence $G(L/F)$ is finite. Let $\mathfrak{p}$ be a prime ideal in $A$. By Lying Over, there exists a prime ideal $\mathfrak{P}$ in $A_L$ lying over $\mathfrak{p}$. We have $(i_L)^{-1}(\mathfrak{p}) = \{\sigma(\mathfrak{P}) \mid \sigma \in G(L/F)\}$ where $i_L : A \to A_L$ is the inclusion map. By Lying Over, if $\mathfrak{Q}$ is a prime ideal in $A_K$, then there exists a prime ideal $\mathfrak{P}'$ in $A_L$ lying over $\mathfrak{Q}$ as $A_L/A_K$ is integral. Hence

$$|\{(\sigma)\}^{-1}(\mathfrak{p})| \leq |\{(\sigma)\}^{-1}(\mathfrak{P})| \leq |G(L/F)| < \infty$$

where $i : A \to A_K$ is the inclusion. □

**Definition 90.19.** A commutative ring $R$ with finitely many maximal ideals is called *semi-local ring*.

**Corollary 90.20.** Let $A$ be a semi-local normal domain with quotient field $F$ and $K/F$ a finite field extension. Then $A_K$ is semi-local.

**Proof.** As $A_K/A$ is integral, $\alpha|_{\text{Max}(A)_{K}} : \text{Max}(A_K) \to \text{Max}(A)$, where $i$ is the inclusion of $A$ in $A_K$. As $\alpha$ has finite fibers, the result follows. □

**Proposition 90.21.** Let $A$ be a normal domain with quotient field $F$ and $K/F$ a finite field extension. Suppose that $\mathfrak{P}$, a prime ideal in $A_K$, lies over the prime ideal $\mathfrak{p}$ in $A$. Let $\overline{F}$ denote the quotient field of $A/\mathfrak{p}$. Then for all $\alpha \in A_K/\mathfrak{P}$, we have $[\overline{F}(\alpha) : \overline{F}] \leq [K : F]$.

**Proof.** Let $\overline{x} : A[t] \to (A/\mathfrak{p})[t] \subset \overline{F}[t]$ be the canonical epimorphism and $x \in A_K$ an element satisfying $\alpha = x + \mathfrak{P}$ in $A_K/\mathfrak{P}$. Then the minimal polynomial $m_F(x)$ of $x$ in $A_K/\mathfrak{P}$ is monic and $\overline{m_F(x)}|_{t = \alpha} = 0$, so

$$[\overline{F}(\alpha) : \overline{F}] \leq \deg \overline{m_F(x)} \leq \deg_F m_F(x) \leq [K : F].$$

□

**Corollary 90.22.** Let $A$ be a normal domain with quotient field $F$ and $K/F$ a finite field extension. Suppose that there exists an element $x$ in $A_K$ satisfying $A_K = A[x]$. Let $\mathfrak{P}$ be a prime ideal in $A_K$ lying over the prime ideal $\mathfrak{p}$ in $A$. Set $\overline{K} = qf(A_K/\mathfrak{P})$ and $\overline{F} = qf(A/\mathfrak{p})$. Then $[\overline{K} : \overline{F}] \leq [K : F]$.

**Proof.** $\overline{K} = qf(A/\mathfrak{p})[x + \mathfrak{P}]$. □

The Cohen-Seidenberg Theorem showed that if $B/A$ is integral, then $\dim A = \dim B$. We are also interested in the heights of primes, as we wish to determine the codimension of irreducible subvarieties. This is more delicate and needs stronger hypotheses. We can now determine one such result.
Theorem 90.23. (Going Down Theorem) Let $B/A$ be integral with $A$ a normal domain and $B$ a domain. Let $p_1 \subset p_2$ be prime ideals in $A$ and $P_2$ a prime ideal in $B$ lying over $p_2$. Then there exists a prime ideal $P_1$ lying over $p_1$ satisfying $P_1 \subset P_2$. In particular, if $P$ is a prime ideal in $B$, then $\text{ht } P = \text{ht}(P \cap A)$.

\[
\begin{array}{c}
\emptyset \subset P_2 \\
| \\
p_1 \subset p_2
\end{array}
\]

Proof. Let $F$ be the quotient field of $A$ and $K$ the quotient field of $B$. Since $B/A$ is integral, $K/F$ is algebraic. Let $L/K$ be a normal closure of $K/F$. As $A_L/A$ is integral and $B \subset A_L$, we have $A_L/B$ is also integral. Let $\mathfrak{O}_1 \in \text{Spec}(A_L)$ lie over $p_1$. By Going Up, there exists a prime ideal $\mathfrak{O}_2 \in V(\mathfrak{O}_1)$ lying over $p_2$. By Lying Over there exists a prime ideal $\mathfrak{O}_2'$ in $A_L$ lying over $\mathfrak{P}_2$ hence over $p_2$. Therefore, there exists an element $\sigma \in G(L/F)$ satisfying $\sigma(\mathfrak{O}_2) = \mathfrak{O}_2'$. We also have $\sigma(\mathfrak{O}_1) \cap A = p_1 = \mathfrak{O}_1 \cap A$. It follows that $\mathfrak{P}_1 = \sigma(\mathfrak{O}_1) \cap B$ works for the first statement.

Let

\[
P_1 < \cdots < P_n = \mathfrak{P}
\]

be a chain of prime ideals in $B$ and $p = \mathfrak{P} \cap A$. By Incomparability, we have

\[
P_1 \cap A < \cdots < P_n \cap A = \mathfrak{P} \cap A = p
\]

a chain of prime ideals in $A$. Therefore, $\text{ht}(\mathfrak{P} \cap A) \geq \text{ht } \mathfrak{P}$. If

\[
p_1 < \cdots < p_n = p = \mathfrak{P} \cap A
\]

is a chain of prime ideals in $A$, then, by Going Down, there exists a chain of prime ideals

\[
P_1 < \cdots < P_n = \mathfrak{P}
\]

in $B$, so $\text{ht } \mathfrak{P} \geq \text{ht}(\mathfrak{P} \cap A)$, and the result follows. \qed

Corollary 90.24. Let $B/A$ be integral with $A$ a normal domain and $B$ a domain. Let $p$ be a prime ideal in $A$ and $\mathfrak{P}$ a prime ideal in $B$ lying over $p$. Then $\text{ht } p = 1$ if and only if $\text{ht } \mathfrak{P} = 1$.

Remark 90.25. The Going Down Theorem 90.23 can be strengthened. In the notation of Theorem 90.23, we can weaken the hypothesis on $B$ to the condition that $B$ be $A$-torsion-free. Indeed, let

\[
S = (A \setminus \{0\})(B \setminus \mathfrak{P}_2) \subset B.
\]

Then $S$ is a multiplicative set excluding $(0)$, since $A$-torsion-free. Hence there exists a prime ideal $\mathfrak{P}$ in $B$ excluding $S$. It follows that $\mathfrak{P}$ satisfies
\( P \subset P_2 \) and \( P \cap A = (0) \). As \( B/A \) induces a ring monomorphism \( A/(P \cap A) \to B/P \), which we view as an inclusion, we are in the situation of the Going Down Theorem. In particular, there exists a prime \( P_1 \) containing \( P \) such that \( P_1/P \) lies over \( p_1 \) and \( P_1/P \subset P_2/P \). Therefore, \( P_1 \) works.

**Exercises 90.26.**

1. Prove the assertions in Remarks 90.3.
2. Prove the assertions in Remarks 90.4.
3. Let \( \varphi : A \to B \) be an integral homomorphism. Show that the associated map \( a_{\varphi} : \text{Spec}(B) \to \text{Spec}(A) \) is a closed map.
4. Prove Corollary 90.12.
5. Let \((R, m)\) be a local ring. Show that any finitely generated \( R \)-projective module (cf. Exercise 38.12(10) and exercises that follow) is a free \( R \)-module.
6. Suppose that \( p_1, \ldots, p_n \) are prime ideals in \( R \) and \( \mathfrak{P} = \bigcap_i p_i \). Show if \( \mathfrak{P} \) is a prime ideal, then there exists an \( i \) with \( \mathfrak{P} = p_i \).
7. Prove Case 2 of Theorem 90.17.
8. (Cayley-Hamilton Theorem) Let \( R \) be a commutative ring, \( \mathfrak{A} \subset R \) an ideal, and \( M \) an \( R \)-module that can be generated by \( n \) elements. If \( \varphi \in \text{End}_R(M) \) satisfies, \( \varphi(M) \subset \mathfrak{A}M \), then there exists a monic polynomial \( f = t^n + a_{n-1}t^{n-1} + \cdots + a_0 \) in \( R[t] \) with coefficients \( a_0, \ldots, a_{n-1} \in \mathfrak{A} \), such that the \( R \)-endomorphism on \( M \), \( f(\varphi) = \varphi^n + a_{n-1}\varphi^{n-1} + \cdots + a_01_M = 0 \).

[The characterization of integral elements (as well as Nakayama’s Lemma 90.10) easy follow from this.]

9. Let \( \mathfrak{A} \) be an ideal in \( R[t] \) and \( - : R[t] \to R[t]/\mathfrak{A} \), the canonical epimorphism. Set \( x = t \) and \( A = F[t]/\mathfrak{A} \). Prove the following:

   (i) As an \( R \)-module \( A \) is generated by at most \( n \) elements if and only if there exists a monic polynomial \( f \in \mathfrak{A} \) of degree at most \( n \). If this is the case, then \( A = \sum_{i=0}^{n-1} Rx_i \), for some \( x_0, \ldots, x_{n-1} \). In particular, \( A \) is a finitely generated \( R \)-module if and only if \( \mathfrak{A} \) contains a monic polynomial.

   (ii) \( A \) is a free \( R \)-module if and only if \( \mathfrak{A} = (f) \) for some monic polynomial \( f \in R[t] \). If this is the case, say rank \( A = n \). Then there exists an \( x \in A \) such that \( \mathcal{B} = \{1, x, \ldots, x^{n-1}\} \) is a basis for \( A \).

10. Under the conditions of Remark 90.25, show if \( p \) is a prime ideal in \( A \) and \( i : A \to B \) is the inclusion map, then

\[
(a_i)^{-1}(p) = \{ \mathfrak{P} \in \text{Spec}(B) \mid \mathfrak{P} \subset V(pB) \text{ is minimal} \}. 
\]
11. Let $R$ be a subring of a commutative ring $A$. Show that the Going Down Theorem holds for $A/R$ if $o_i : \text{Spec}(A) \to \text{Spec}(R)$ is an open map. [The converse of this is true if $R$ is a Noetherian ring.]

91. Primary Decomposition

As throughout this chapter, $R$ will denote a commutative ring. In this section, we investigate the generalization of unique factorization of ideals that Dedekind domains satisfy called the Lasker-Noether Theorem. As unique factorization of ideals into a product of prime ideals characterize Dedekind domains, the appropriate generalization is, as one would expect, weaker. Multiplication of ideals is replaced by intersections, powers of primes by primary ideals, and not all primary ideals are unique. However, this generalization will apply to any Noetherian ring, not just a Noetherian domain.

We shall need the following lemma whose proof we leave as an exercise:

Recall that the radical of an ideal $\mathfrak{a}$ in $R$ is defined to be

\[ \sqrt{\mathfrak{a}} := \{ x \in R \mid x^n \in \mathfrak{a} \text{ for some } n \in \mathbb{Z}^+ \}. \]

**Lemma 91.1.** Let $\mathfrak{a}_1, \ldots, \mathfrak{a}_n$ be ideals in $R$. Then

\[ \sqrt{\bigcap_{i=1}^n \mathfrak{a}_i} = \bigcap_{i=1}^n \sqrt{\mathfrak{a}_i} = \sqrt{\mathfrak{a}_1 \cdots \mathfrak{a}_n}. \]

Recall that an ideal $\mathfrak{q} < R$ is is called a primary ideal if it satisfies one of the following equivalent conditions:

(i) If $xy$ lies in $\mathfrak{q}$ then either $x$ lies in $\mathfrak{q}$ or there exists a positive integer $n$ with $y^n \in \mathfrak{q}$.

(ii) If $xy$ lies in $\mathfrak{q}$ then either $x$ lies in $\mathfrak{q}$ or $y \in \sqrt{\mathfrak{q}}$.

(iii) The set of zero divisors $\text{zd}(R/\mathfrak{a})$ of $R/\mathfrak{a}$ lies in the nilradical $\text{nil}(R/\mathfrak{q})$ of $R/\mathfrak{a}$.

**Lemma 91.2.** Let $\mathfrak{q}$ be a primary ideal in $R$. Then $\sqrt{\mathfrak{q}}$ is a prime ideal in $R$ and is the smallest prime ideal in $R$ containing $\mathfrak{q}$.

**Proof.** Suppose that $xy$ lies in $\mathfrak{q}$. Then there exists a positive integer $n$ such that $x^n y^n$ lies in $\mathfrak{q}$. Hence either $x^n \in \mathfrak{q}$ or $y^n \in \mathfrak{q}$ i.e., either $x \in \sqrt{\mathfrak{q}}$ or $y \in \sqrt{\mathfrak{q}}$. It follows that $\sqrt{\mathfrak{q}}$ is a prime ideal in $R$. As $\mathfrak{q} \subseteq \sqrt{\mathfrak{q}}$ and if $\mathfrak{q}$ lies in a prime ideal $\mathfrak{p}$ so does $\sqrt{\mathfrak{q}}$, the second statement also follows. \(\square\)

**Definition 91.3.** If $\mathfrak{q}$ be a primary ideal in $R$ and $\mathfrak{p}$ is the prime ideal $\sqrt{\mathfrak{q}}$, we say that $\mathfrak{q}$ is a $\mathfrak{p}$-primary ideal.
A useful example is given by the following result:

**Lemma 91.4.** Let \( \mathfrak{A} < R \) be an ideal satisfying \( \sqrt{\mathfrak{A}} \) is a maximal ideal in \( R \). Then \( \mathfrak{A} \) is a \( \sqrt{\mathfrak{A}} \)-primary ideal. In particular, if \( \mathfrak{m} \) is a maximal ideal in \( R \), then \( \mathfrak{m}^n \) is an \( \mathfrak{m} \)-primary ideal for each positive integer \( n \).

**Proof.** Let \( \mathfrak{m} = \sqrt{\mathfrak{A}} \). Then \( \mathfrak{m} = \bigcap_{V(\mathfrak{A})} \mathfrak{p} \), hence \( \mathfrak{m} \subset \mathfrak{p} \) for every prime ideal in \( V(\mathfrak{A}) \). Since \( \mathfrak{m} \) is maximal, \( V(\mathfrak{A}) = \{ \mathfrak{m} \} \). It follows that \( R/\mathfrak{A} \) is a primary ring, i.e., \( |\text{Spec}(R/\mathfrak{A})| = 1 \). In particular, we have \( \text{zd}(R/\mathfrak{A}) = \text{nil}(R/\mathfrak{A}) \), so \( \mathfrak{A} \) is \( \mathfrak{m} \)-primary.

**Examples 91.5.**

1. Each prime ideal in \( R \) is a primary ideal in \( R \).
2. The ideal \((n)\) in \( \mathbb{Z} \) is primary if and only if \( n = 0 \) or \( n \) is a power of a prime.
3. Let \( F \) be a field and \( x, y, z \) indeterminants over \( F \). Then \((xy - z^2)\) is a prime ideal in the UFD \( F[x, y, z] \). Set \( R = F[x, y, z]/(xy - z^2) \) and \( \phi : F[x, y, z] \to R = F[\bar{x}, \bar{y}, \bar{z}] \) the canonical epimorphism. Then \( \mathfrak{p} = (\bar{z}) \) is a prime ideal in the domain \( R \). We have \( \pi y = z^2 \in p^2 := \mathfrak{p}\mathfrak{p} \), but \( \pi^2 \notin p^2 \) and \( \bar{y} \notin \sqrt{p^2} = \mathfrak{p} \).

Therefore, \( p^2 \) is not a primary ideal in \( R \). So, in general, a power of a prime ideal may not be a primary ideal.
4. Let \( F \) be a field, \( x, y \) indeterminants over \( F \) and \( R = F[x, y] \). Then \( \mathfrak{Q} = (x, y^2) \) is an \( (x, y) \)-primary ideal in \( R \) as \( R/\mathfrak{Q} = F[y]/(y^2) \) satisfies \( \text{zd}(R/\mathfrak{Q}) \subset \text{nil}(R/\mathfrak{Q}) \). Let \( \mathfrak{p} = (x, y) \), then \( (x^2, xy, y^2) = p^2 < \mathfrak{Q} \subset \mathfrak{p} \). It follows that \( \mathfrak{Q} \neq p^n \) for any integer \( n \), i.e., a primary ideal may not be a power of a prime ideal.
5. Suppose that \( \mathfrak{Q}_i, i = 1, \ldots, r \), are all \( \mathfrak{p} \)-primary ideals. Then so is \( \bigcap \mathfrak{Q}_i \):

We know that \( \sqrt{\bigcap \mathfrak{Q}_i} = \bigcap \sqrt{\mathfrak{Q}_i} = \mathfrak{p} \). If \( xy \in \bigcap \mathfrak{Q}_i \) with \( y \notin \bigcap \mathfrak{Q}_i \), then there exists a \( j \) with \( y \notin \mathfrak{Q}_j \). As \( xy \in \mathfrak{Q}_j \), we have \( x^n \) lies in \( \mathfrak{Q}_j \) for some \( n \), hence \( x \) lies in \( \sqrt{\mathfrak{Q}_j} = \mathfrak{p} = \sqrt{\bigcap \mathfrak{Q}_i} \), so \( \bigcap \mathfrak{Q}_i \) is \( \mathfrak{p} \)-primary.

**Definition 91.6.** Let \( \mathfrak{A} \) and \( \mathfrak{B} \) be two ideals in \( R \). Set

\[(\mathfrak{A} : \mathfrak{B}) := \{ y \in R \mid y\mathfrak{B} \subset \mathfrak{A} \} \subset R,\]

an ideal in \( R \) called the colon ideal.

If \( \mathfrak{B} = (x) \), we write \( (\mathfrak{A} : x) \) for \( (\mathfrak{A} : (x)) \). If \( \mathfrak{A} = 0 \), then \( (0 : x) = \{ y \in R \mid yx = 0 \} \) is the annihilator \( \text{ann}_R(x) \) of \( x \) in \( R \).

Note the following:

**Remarks 91.7.** We have:
1. \( \text{zd}(R) = \bigcup_{x \neq 0} \sqrt{\text{ann}_R(x)} = \bigcup_{x \neq 0} \sqrt{0 : x} \).

2. If \( \mathfrak{A}_i \subset R \) are ideals for \( i = 1, \ldots, n \), then
\[
(\bigcap_{i=1}^{n} \mathfrak{A}_i : x) = \bigcap_{i=1}^{n} (\mathfrak{A}_i : x).
\]

**Computation 91.8.** Let \( \mathfrak{Q} \) be a \( p \)-primary ideal in \( R \) and \( x \in R \). Then we have:

(i) If \( x \in \mathfrak{Q} \), then \( (\mathfrak{Q} : x) = R \).

(ii) If \( x \notin \mathfrak{Q} \), then \( (\mathfrak{Q} : x) \) is \( p \)-primary. In particular, \( \sqrt{(\mathfrak{Q} : x)} = p \).

(iii) If \( x \notin p \), then \( (\mathfrak{Q} : x) = \mathfrak{Q} \).

**Proof.** (i) is immediate.

(iii): If \( x \notin p \) and \( y \in (\mathfrak{Q} : x) \), then \( xy \in \mathfrak{Q} \), hence \( y \in \mathfrak{Q} \) by definition.

(ii): Let \( y \in (\mathfrak{Q} : x) \). As \( x \notin p \) and \( xy \in \mathfrak{Q} \), we have \( y \in \sqrt{\mathfrak{Q}} = p \).

Hence \( \mathfrak{Q} \subset (\mathfrak{Q} : x) \subset p \), so \( p = \sqrt{\mathfrak{Q}} \subset \sqrt{(\mathfrak{Q} : x)} \subset \sqrt{p} = p \).

Next let \( y, z \in R \) with \( yz \in (\mathfrak{Q} : x) \) and \( y \notin \sqrt{\mathfrak{Q}} = p \). Then \( yzx \in \mathfrak{Q} \) implies that \( xz \in \mathfrak{Q} \), hence \( z \in (\mathfrak{Q} : x) \) showing that \( (\mathfrak{Q} : x) \) is \( p \)-primary. \( \square \)

With the above definitions, we can now define the type of decomposition of ideals in which we shall be interested.

**Definition 91.9.** Let \( \mathfrak{A} < R \) be an ideal. A primary decomposition of \( \mathfrak{A} \) is an equation
\[
(*) \quad \mathfrak{A} = \mathfrak{Q}_1 \cap \cdots \cap \mathfrak{Q}_n \text{ with } \mathfrak{Q}_i \text{ primary for } i = 1, \ldots n.
\]

We say that \( (*) \) is irredundant if, in addition,

(i) \( \sqrt{\mathfrak{Q}_1}, \ldots, \sqrt{\mathfrak{Q}_n} \) are all distinct.

(ii) \( \bigcap_{j \neq i} \mathfrak{Q}_j \not\subset \mathfrak{Q}_i \) for each \( i = 1, \ldots, n \).

**Remark 91.10.** Suppose that an ideal \( \mathfrak{A} \) has a primary decomposition \( (*) \). By Example 91.5(5), if \( \mathfrak{Q}_i, \ldots, \mathfrak{Q}_s \) have the same radical, then \( \mathfrak{Q}_i \cap \cdots \cap \mathfrak{Q}_s \) is \( \sqrt{\mathfrak{Q}_i} \)-primary, so any primary decomposition \( (*) \) can be made to satisfy (i). Throwing out those \( \mathfrak{Q}_i \)'s not satisfying (ii), then shows that any \( \mathfrak{A} \) having a primary decomposition has an irredundant primary decomposition.

We show our first uniqueness statement for irredundant primary decompositions.
Theorem 91.11. Suppose that \( \mathfrak{A} \subset R \) is an ideal that has an irredundant primary decomposition
\[
\mathfrak{A} = \mathfrak{Q}_1 \cap \cdots \cap \mathfrak{Q}_n \quad \text{with} \quad p_i = \sqrt{\mathfrak{Q}_i} \quad \text{for} \quad i = 1, \ldots, n.
\]
Then
\[
\{ p_1, \ldots, p_n \} = \{ \sqrt{(\mathfrak{A} : x)} \mid x \in R \} \cap \text{Spec}(R)
\]
and is independent of the irredundant primary decomposition, i.e., the radicals of the primary ideals giving any irredundant primary decomposition of \( \mathfrak{A} \) are unique.

Proof. Let \( x \in R \). We have \( (\mathfrak{A} : x) = (\cap \mathfrak{Q}_i : x) = \cap (\mathfrak{Q}_i : x) \), so by the Computation 91.8,
\[
\sqrt{(\mathfrak{A} : x)} = \bigcap \sqrt{(\mathfrak{Q}_i : x)} = \bigcap_{x \not\in \mathfrak{Q}_i} p_i.
\]
If \( \sqrt{(\mathfrak{A} : x)} \) is a prime ideal, it follows that there exists a \( \mathfrak{Q}_i \) with \( x \not\in \mathfrak{Q}_i \) such that \( \sqrt{(\mathfrak{A} : x)} = p_i \) by Exercise 90.26(6).

Conversely, since our primary decomposition is irredundant, for all \( i \) and all \( x_i \) satisfying \( x_i \in \bigcap_{j \neq i} \mathfrak{Q}_j \setminus \mathfrak{Q}_i \), we have \( \sqrt{(\mathfrak{A} : x_i)} = p_i \) by Computation 91.8. \( \square \)

The proof of the above theorem and Computation 91.8 show

Corollary 91.12. If \( \mathfrak{A} \subset R \) is an ideal having an irredundant primary decomposition \( \mathfrak{A} = \mathfrak{Q}_1 \cap \cdots \cap \mathfrak{Q}_n \) with \( p_i = \sqrt{\mathfrak{Q}_i} \) for \( i = 1, \ldots, n \), then there exists an element \( x_i \) in \( \mathfrak{Q}_i \) such that \( (\mathfrak{A} : x_i) \) is \( p_i \)-primary.

Definition 91.13. Let \( \mathfrak{A} = \mathfrak{Q}_1 \cap \cdots \cap \mathfrak{Q}_n \) be an irredundant primary decomposition of \( \mathfrak{A} \) in \( R \) with \( p_i = \sqrt{\mathfrak{Q}_i} \) for \( i = 1, \ldots, n \). We set
\[
\text{Ass}_R V(\mathfrak{A}) := \{ p_1, \ldots, p_n \},
\]
the set of associated prime ideals of \( \mathfrak{A} \) [actually of \( R/\mathfrak{A} \) in modern language]. We partially order \( \text{Ass}_R V(\mathfrak{A}) \) by set inclusion \( \subset \). Elements that are minimal in \( \text{Ass}_R V(\mathfrak{A}) \) under this partial order are called isolated prime ideals of \( \mathfrak{A} \). The others (if any) are called embedded prime ideals of \( \mathfrak{A} \).

Isolated prime ideals in \( \text{Ass}_R V(\mathfrak{A}) \) are, in fact, precisely the prime ideals minimally containing \( \mathfrak{A} \) in \( R \), which we shall show next. In the sequel we shall write \( p \in V(\mathfrak{A}) \) is minimal for such a prime ideal.

Proposition 91.14. Suppose that the ideal \( \mathfrak{A} \) has an irredundant primary decomposition. Then a prime ideal \( p \) in \( R \) is an isolated prime of \( \mathfrak{A} \) if and only if \( p \in V(\mathfrak{A}) \) is minimal and only if \( p/\mathfrak{A} \in \text{Min}(R/\mathfrak{A}) \).
Proof. Let $\mathfrak{A} = \mathfrak{Q}_1 \cap \cdots \cap \mathfrak{Q}_n$ be an irredundant primary decomposition of $\mathfrak{A}$ in $R$ with $p_i = \sqrt{\mathfrak{Q}_i}$ for $i = 1, \ldots, n$. Suppose that $p$ lies in $V(\mathfrak{A})$. Then
\[ p = \sqrt{p} \supseteq \sqrt{\mathfrak{A}} = \cap \sqrt{\mathfrak{Q}_i} = \cap p_i. \]
Hence there exists an $i$ such that $p_i \subset p$, i.e., $p$ contains an isolated prime of $\mathfrak{A}$. \qed

Because of the proposition, isolated primes are also called minimal primes of $\mathfrak{A}$.

Proposition 91.15. Suppose that the ideal $\mathfrak{A}$ in $R$ has an irredundant primary decomposition $\mathfrak{A} = \mathfrak{Q}_1 \cap \cdots \cap \mathfrak{Q}_n$ with $p_i = \sqrt{\mathfrak{Q}_i}$ for $i = 1, \ldots, n$. Then
\[ \bigcup p_i = \{ x \in R \mid (\mathfrak{A} : x) > \mathfrak{A} \}. \]
In particular, if $\mathfrak{A} = (0)$, then
\[ zd(R) = \bigcup_{\text{Ass}_R(0)} p. \]

Proof. Let $- : R \to R/\mathfrak{A}$ be the canonical epimorphism.

Check. $(0) = \overline{\mathfrak{Q}_1 \cap \cdots \cap \mathfrak{Q}_n}$ is an irredundant primary decomposition in $\overline{R}$. So it suffices to prove $zd(R) = \bigcup_{\text{Ass}_R(0)} p$.

By the proof of Theorem 91.11, each $p_j = \sqrt{(0 : x)}$ for some $x \in R$. The result follows. \qed

So we now know if $(0)$ has an irredundant primary decomposition, then we have
\[ zd(R) = \bigcup_{\text{Ass}_R(0)} p \]
\[ \text{nil}(R) = \bigcap_{\text{Spec}(R)} p = \bigcap_{\text{Min}(R)} p \]
\[ \text{Min}(R) \subset \text{Ass}_R V(0). \]

Remark 91.16. If they exist, irredundant primary decompositions are not necessarily unique. An example is given as follows: Let $R = F[x, y]$ with $F$ a field and $x, y$ indeterminants. Then
\[ (x) \cap (x, y)^2 = (x^2, xy) = (x) \cap (x^2, y) \]
are two irredundant primary decompositions for $(x^2, xy)$.

We shall obtain a ‘weaker’ uniqueness statement. We leave the proof of the following as an exercise.
Lemma 91.17. Primary ideals satisfy:

(1) Let \( S \) be a multiplicative set in \( R \) and \( \varphi_R : R \to S^{-1}R \) the canonical ring homomorphism given by \( r \mapsto r/1 \). Then \( \varphi_R \) induces a bijection:
\[
\{ Q \mid Q \subset R \text{ primary with } Q \cap S = \emptyset \} \quad \mapsto \quad \{ Q \mid Q \subset S^{-1}R \text{ primary} \}
\]
via \( Q \mapsto S^{-1}Q \).

(2) If \( \psi : A \to B \) is a ring homomorphism of commutative rings, \( Q \subset B \) a primary ideal, then \( \psi^{-1}(Q) \subset A \) is a primary ideal.

Proposition 91.18. Suppose that the ideal \( \mathfrak{A} \) in \( R \) has an irredundant primary decomposition \( \mathfrak{A} = \mathfrak{Q}_1 \cap \cdots \cap \mathfrak{Q}_m \) with \( \sqrt{\mathfrak{Q}_i} = p_i \) for \( i = 1, \ldots, m \) and \( S \) is a multiplicative set in \( R \) satisfying \( S \cap \mathfrak{Q}_i = \emptyset \) for \( 1 \leq i \leq m \), and \( S \cap \mathfrak{Q}_i \neq \emptyset \) for \( i > m \). Then
\[
S^{-1}\mathfrak{A} = S^{-1}\mathfrak{Q}_1 \cap \cdots \cap S^{-1}\mathfrak{Q}_m
\]
is an irredundant primary decomposition for \( S^{-1}\mathfrak{A} \).

Proof. We have
\[
S^{-1}\mathfrak{A} = S^{-1}(\bigcap_{i=1}^{m} \mathfrak{Q}_i) = \bigcap_{i=1}^{m} S^{-1}\mathfrak{Q}_i = \bigcap_{i=1}^{m} S^{-1}\mathfrak{Q}_i
\]
with \( S^{-1}\mathfrak{Q}_i \) being \( S^{-1}p_i \)-primary, \( i = 1, \ldots, m \) and the \( S^{-1}p_1, \ldots, S^{-1}p_m \) distinct. It follows that \( S^{-1}\mathfrak{A} = \bigcap_{i=1}^{m} S^{-1}\mathfrak{Q}_i \) is an irredundant primary decomposition. \( \square \)

Remark 91.19. In the above, if \( \varphi_R : R \to S^{-1}R \) is the canonical ring homomorphism \( r \mapsto r/1 \), then
\[
(\varphi_R)^{-1}(S^{-1}\mathfrak{A}) = (\varphi_R)^{-1}(\bigcap_{i=1}^{m} S^{-1}\mathfrak{Q}_i) = \bigcap_{i=1}^{m} (\varphi_R)^{-1}(S^{-1}\mathfrak{Q}_i) = \bigcap_{i=1}^{m} \mathfrak{Q}_i.
\]

Definition 91.20. Let \( \mathfrak{A} \subset R \) be an ideal with an irredundant primary decomposition. A subset \( I \subset \text{Ass}_RV(\mathfrak{A}) \) is called isolated if whenever \( p, p' \) lie in \( \text{Ass}_RV(\mathfrak{A}) \) and satisfy \( p \subset p' \), then \( p' \in I \) implies \( p \in I \), i.e., all descending chains of prime ideals lying in \( \text{Ass}_RV(\mathfrak{A}) \) beginning with an element of \( I \) have all elements in the chain lying in \( I \).

Remark 91.21. Let \( I \subset \text{Ass}_RV(\mathfrak{A}) \) be an isolated set where \( \mathfrak{A} \) has an irredundant primary decomposition. Let \( S = R \setminus \bigcup_i p_i \). Then \( S \) is a saturated multiplicative set. Moreover, if \( p' \in \text{Ass}_RV(\mathfrak{A}) \), then we have:

(i) If \( p' \in I \), then \( p' \cap S = \emptyset \).
(ii) If $p' \notin I$, then $p' \not\subseteq \bigcup_{j} p$ (by the Prime Avoidance Lemma 90.16), hence $p' \cap S \neq \emptyset$.

In particular, we obtain our second uniqueness result:

**Theorem 91.22.** Suppose that the ideal $\mathfrak{A}$ in $R$ has an irredundant primary decomposition $\mathfrak{A} = \mathfrak{Q}_1 \cap \cdots \cap \mathfrak{Q}_n$ with $p_i = \sqrt{\mathfrak{Q}_i}$ for $i = 1, \ldots, n$ and $I = \{p_1, \ldots, p_m\} \subset \text{Ass}_R V(\mathfrak{A})$ an isolated subset. Then $\bigcap_{i=1}^m \mathfrak{Q}_i$ is independent of an irredundant primary decomposition for $\mathfrak{A}$. In particular, the $\mathfrak{Q}_i$ with $p_i \in V(\mathfrak{A})$ minimal are uniquely determined by $\mathfrak{A}$, i.e., the isolated primes determine unique primary ideals.

**Proof.** Let $S = R \setminus \bigcup_{j=1}^m p_j$. Then by the above remark,

$$\mathfrak{Q}_i \cap \cdots \cap \mathfrak{Q}_n = \varphi^{-1}_R(S^{-1}\mathfrak{A})$$

where $\varphi_R : R \rightarrow S^{-1}R$ is the canonical ring homomorphism $r \mapsto r/1$.

**Example 91.23.** Let $\mathfrak{A} < R$ be a radical ideal, i.e., $\mathfrak{A} = \sqrt{\mathfrak{A}}$. Suppose that $|\text{Min}(R/\mathfrak{A})|$ is finite. Then $\mathfrak{A}$ has an irredundant primary decomposition and $\text{Ass}_R V(\mathfrak{A}) = \{p \in V(\mathfrak{A}) \mid p \text{ minimal}\}$, i.e., there are no embedded primes:

We have

$$\mathfrak{A} = \sqrt{\mathfrak{A}} = \bigcap_{p \in V(\mathfrak{A}) \text{ minimal}} p$$

with the right hand side an irredundant primary decomposition of $\mathfrak{A}$. If $R$ is Noetherian, we have established in Corollary 89.27 and will establish again below that $|\text{Min}(R/\mathfrak{A})|$ is always finite.

We wish to show that every ideal in a Noetherian ring has an irredundant primary decomposition. To do this we shall use Exercise 29.22(15) which asked to show

**Lemma 91.24.** An ideal $\mathfrak{C}$ in a commutative ring $R$ is called irreducible if whenever $\mathfrak{C} = \mathfrak{A} \cap \mathfrak{B}$ for some ideals $\mathfrak{A}$ and $\mathfrak{B}$ in $R$, then either $\mathfrak{C} = \mathfrak{A}$ or $\mathfrak{C} = \mathfrak{B}$. If $R$ is Noetherian, then every ideal $\mathfrak{A} < R$ is a finite intersection of irreducible ideals of $R$, i.e., $\mathfrak{A} = \mathfrak{C}_1 \cap \cdots \cap \mathfrak{C}_n$, for some irreducible ideals $\mathfrak{C}_i$ in $R$. We call such a decomposition an irreducible decomposition of $\mathfrak{A}$.

This is what is needed as:

**Lemma 91.25.** Let $R$ be Noetherian ring and $\mathfrak{A} < R$ an irreducible ideal. Then $\mathfrak{A}$ is a primary ideal.
Proof. Since $R/\mathfrak{A}$ is also Noetherian, by the Correspondence Theorem, it suffices to show that if $(0)$ is irreducible, then it is primary. So suppose that $(0)$ is irreducible and $xy = 0$, $x, y \in R$, with $y \neq 0$. By the ascending chain condition,

$$
(0 : x) \subset (0 : x^2) \subset \cdots \subset (0 : x^n) \subset \cdots
$$

stabilizes, say $(0 : x^N) = (0 : x^{N+i})$ for all $i \geq 1$.

Claim. $(0) = (y) \cap (x^N)$.

If we prove the claim, the result will follow as $(0)$ being irreducible implies $x^N = 0$ by the exercise. To show the claim, let $z \in (y) \cap (x^N)$. As $z \in (y)$, we have $xz = 0$ and as $z \in (x^N)$, we have $z = ax^N$ for some $a \in R$. Consequently, $0 = xz = ax^{N+1}$, so $a \in (0 : x^{N+1}) = (0 : x^N)$. Therefore, $z = ax^N = 0$ as needed. □

Putting all of what we have done together establishes:

Theorem 91.26. (Lasker-Noether Theorem) Let $R$ be a commutative Noetherian ring and $\mathfrak{A} < R$ an ideal. Then $\mathfrak{A}$ has an irredundant primary decomposition, say

$$
\mathfrak{A} = \mathfrak{Q}_1 \cap \cdots \cap \mathfrak{Q}_n
$$

with $p_i = \sqrt{\mathfrak{Q}_i}$ for $i = 1, \ldots, n$.

Moreover,

1. $p_1, \ldots, p_n$ are unique.
2. $\mathfrak{Q}_i$ is unique if $p_i$ is isolated.

In particular, $|\text{Min}(R/\mathfrak{A})| \leq |\text{Ass}_R \mathfrak{A}|$ is finite.

We derive some consequences.

Corollary 91.27. Let $R$ be a Noetherian ring. Then $\text{Min}(R)$ is finite.

Lemma 91.28. Let $R$ be a Noetherian ring and $\mathfrak{A} < R$ an ideal. Then there exists a positive integer $n$ satisfying $(\sqrt{\mathfrak{A}})^n \subset \mathfrak{A}$.

Proof. Since ideals in a Noetherian ring are finitely generated, $\mathfrak{A} = (a_1, \ldots, a_n)$ for some $a_i \in \sqrt{\mathfrak{A}}$. It follows that there exists a positive integer $N$ such that $a_i^N \in \mathfrak{A}$ for all $i$, so $(\sqrt{\mathfrak{A}})^N \subset \mathfrak{A}$. □

Corollary 91.29. Let $R$ be a Noetherian ring. Then $\text{nil}(R)$ is nilpotent, i.e., there exists a positive integer $N$ such that $(\text{nil}(R))^N = 0$.

Corollary 91.30. Let $R$ be a Noetherian ring, $\mathfrak{m} \in \text{Max}(R)$, and $\mathfrak{Q} < R$ an ideal. Then the following are equivalent:

1. $\mathfrak{Q}$ is $\mathfrak{m}$-primary.
2. $\sqrt{\mathfrak{Q}} = \mathfrak{m}$.
3. There exists a positive integer $n$ such that $\mathfrak{m}^n \subset \mathfrak{Q} \subset \mathfrak{m}$. 


Proof. We have shown all but the implication (3) \(\Rightarrow\) (2) which follows from \(m = \sqrt{m^n} = \sqrt{m} = m\). \(\square\)

**Proposition 91.31.** Let \(R\) be a Noetherian ring and \(\mathfrak{A} < R\) an ideal. Then

\[
\text{Ass}_R V(\mathfrak{A}) = \{(\mathfrak{A}: x) \mid x \in R\} \cap \text{Spec}(R)
\]

**Proof.** Replacing \(R\) with \(R/\mathfrak{A}\), we may assume that \(\mathfrak{A} = 0\). Let \(\mathfrak{A} = \mathfrak{Q}_1 \cap \cdots \cap \mathfrak{Q}_n\) be an irredundant primary decomposition with \(\mathfrak{p}_i = \sqrt{\mathfrak{Q}_i}\) for \(i = 1, \ldots, n\). As the decomposition is irredundant, \(\mathfrak{A}_i = \cap_{j \neq i} \mathfrak{Q}_j > 0\). By the proof of Theorem 91.11, \(\mathfrak{p}_i = \sqrt{(0 : x)}\) for all nonzero \(x\) in \(\mathfrak{A}_i\), since \(\mathfrak{A}_i \cap \mathfrak{Q}_i = 0\). Hence

\[
(0 : x) = \text{ann}_R x \subset \mathfrak{p}_i \text{ for all nonzero } x \in \mathfrak{A}_i.
\]

As \(\mathfrak{Q}_i\) is \(\mathfrak{p}_i\)-primary, there exists a positive integer \(m\) satisfying \(\mathfrak{p}_i^m \subset \mathfrak{Q}_i\) by Lemma 91.28. Consequently, \(\mathfrak{A}_i \mathfrak{p}_i^m \subset \mathfrak{A}_i \cap \mathfrak{p}_i^m \subset \mathfrak{A}_i \cap \mathfrak{Q}_i = 0\).

Choose \(m > 0\) minimal with \(\mathfrak{A}_i \cap \mathfrak{p}_i^m = 0\). If \(x \in \mathfrak{A}_i \mathfrak{p}_i^{m-1} \subset \mathfrak{A}_i\), then \(xp_i = 0\), so we conclude that \(\mathfrak{p}_i \subset \text{ann}_R x = (0 : x)\) if \(x\) is nonzero.

Conversely, if \(\mathfrak{p} = (0 : x)\) is a prime ideal, then \(\sqrt{(0 : x)} = \sqrt{\text{ann}_R x} = \mathfrak{p}\), and \(\mathfrak{p} \in \text{Ass}_R V(0)\). \(\square\)

**Exercises 91.32.**
1. Prove Lemma 91.1
2. Prove Lemma 91.4
3. Show that the colon ideal \((\mathfrak{A} : \mathfrak{B})\) of the ideals \(\mathfrak{A}\) and \(\mathfrak{B}\) in \(R\) is an ideal in \(R\) and the largest ideal \(\mathfrak{C}\) in \(R\) satisfying \(\mathfrak{BC} \subset \mathfrak{A}\).
4. Prove Remark 91.7
5. Prove Lemma 91.17

**92. Addendum: Associated Primes of Modules**

In the Section §91, we gave the best generalization of the Fundamental Theorem of Algebra, thereby concluding this part of our investigations into analogues of unique factorization domains. Some of the constructions and ideas in that section have become more important than primary decomposition itself. Many are easier to develop than in that section, so we generalize the approach and get alternative proofs of some of the results in Section 91. In particular, if \(\mathfrak{A}\) is an ideal in a commutative ring \(R\), then the set of associated primes of \(\mathfrak{A}\) were denoted \(\text{Ass}_R V(\mathfrak{A})\). In this section, this will become \(\text{Ass}_R (R/\mathfrak{A})\) as we shall define the associated primes of an \(R\)-module. Throughout this section \(R\) is a commutative ring.
**Definition 92.1.** Let \( R \) be a commutative ring and \( M \) an \( R \)-module. A prime ideal \( p \) in \( R \) is called an associated prime of \( M \) if there exists a nonzero element \( m \) in \( M \) such that \( p = \text{ann}_R m \). Let

\[
\text{Ass}_R(M) := \{ p \in \text{Spec} R \mid p \text{ is an associated prime of } M \}.
\]

**Remarks 92.2.** Let \( M \) be an \( R \)-module and \( p \) a prime ideal in \( R \).

1. The prime ideal \( p \) lies in \( \text{Ass}_R(M) \) if and only if there exists an \( R \)-monomorphism \( R/p \to M \), since \( R/\text{ann}_R m \cong Rm \) for all \( m \) in \( M \).
2. \( R/p \) is a domain, so \( p = \text{ann}_R(r + p) \) for all \( r \in R \setminus p \), hence

\[
\text{Ass}_R(R/p) = \{ p \}.
\]

**Lemma 92.3.** Let \( R \) be a commutative ring and

\[
0 \to M' \to M \to M'' \to 0
\]

be an exact sequence of \( R \)-modules. Then

\[
\text{Ass}_R(M') \subset \text{Ass}_R(M) \subset \text{Ass}_R(M') \cup \text{Ass}_R(M'').
\]

**Proof.** In the exact sequence, we may assume that \( M' \subset M \), \( M'' = M/M' \), and \( \tilde{\cdot} : M \to M/M' \) is the canonical \( R \)-homomorphism. If \( p \in \text{Ass}_R(M') \), there exists an \( R \)-monomorphism \( R/p \to M' \subset M \), so \( p \in \text{Ass}_R(M) \).

Suppose that \( p \in \text{Ass}_R(M) \). Let \( N \) be the image of the \( R \)-monomorphism \( R/p \to M \). If \( N \cap M' = 0 \), then the restriction of \( \tilde{\cdot} : M \to M/M' \) to \( N \) gives an \( R \)-monomorphism \( N \to M/M' \), so \( p \in \text{Ass}_R(M/M') \). If \( 0 \neq N \cap M' \subset N \), since \( p = \text{ann}_R n \) for all nonzero \( n \) in \( N \), we have \( p \in \text{Ass}_R(N \cap M') \subset \text{Ass}_R(M') \). □

We shall see that the second set inclusion in the lemma can be proper. [The first set inclusion can easily seen to be proper in general.]

The generalization of Proposition 91.31 is easy to prove.

**Proposition 92.4.** Let \( R \) be a commutative ring and \( M \) be a nonzero \( R \)-module. Suppose that

\[
\mathfrak{A} \in \{ \text{ann}_R m \mid 0 \neq m \in M \} \text{ is maximal relative to } \subset.
\]

Then \( \mathfrak{A} \in \text{Ass}_R(M) \). In particular if \( R \) is Noetherian and \( N \) is an \( R \)-module, then

\[
\text{Ass}_R(N) = \emptyset \text{ if and only if } N = 0.
\]

**Proof.** Suppose that there exist \( a \) and \( b \) in \( R \) with \( b \notin \mathfrak{A} \) but \( ab \in \mathfrak{A} \). By definition \( \mathfrak{A} = \text{ann}_R m \) for some nonzero element \( m \) in \( M \). As \( bm \neq 0 \), we have \( \text{ann}_R m \subset \text{ann}_R bm < R \). By maximality,
ann\(_R\) m = ann\(_R\) bm, so abm = 0 implies that \(a \in \text{ann}_R bm \subseteq \text{ann}_R m = \mathfrak{A}\). Therefore, \(\mathfrak{A} \in \text{Spec}(R)\). Furthermore, if \(R\) is Noetherian and \(N\) is nonzero, then \(\{\text{ann}_R m \mid 0 \neq n \in N\}\) has a maximal element. \(\square\)

Using the proposition, we get a simple proof to the generalization of the conclusion of Proposition 91.15.

**Corollary 92.5.** Let \(R\) be a Noetherian ring and \(M\) a nonzero \(R\)-module. Then \(zd(M) = \bigcup_{\text{Ass}_R(M)} \mathfrak{p}\).

**Proposition 92.6.** Let \(R\) be a Noetherian ring and \(M\) a nonzero finitely generated \(R\)-module. Then there exists a chain of \(R\)-modules

\[
0 = M_0 < M_1 < \cdots < M_n = M
\]

for some \(n\) satisfying

\[
M_i/M_{i-1} \cong R/\mathfrak{p}_i \quad \text{for some prime ideals} \quad \mathfrak{p}_i, \; i = 1, \ldots, n.
\]

**Proof.** Let \(\mathcal{S} = \{N \mid 0 < N \subset M \text{ } R\text{-modules with } N \text{ satisfying the Proposition}\}\).

Since \(M \neq 0\), there exists a prime ideal \(\mathfrak{p} \in \text{Ass}_R(M)\) and an \(R\)-monomorphism \(\varphi : R/\mathfrak{p} \hookrightarrow M\). Therefore, \(\text{im}\varphi\) lies in \(\mathcal{S}\), i.e., \(\mathcal{S}\) is non-empty. As \(R\) is Noetherian and \(M\) is finitely generated, \(M\) is a Noetherian \(R\)-module. Hence there there exists a maximal element \(N \in \mathcal{S}\) (relative to \(\subset\)). Suppose that \(N < M\). Then there exists a prime ideal in \(\text{Ass}_R(M/N)\). Consequently, there exists an \(R\)-module \(N'\) satisfying \(N < N' \subset M\) and an \(R\)-isomorphism \(R/\mathfrak{p} \rightarrow N'/N\). As \(N \in \mathcal{S}\) and \(N'/N \cong R/\mathfrak{p}\), we have \(N' \in \mathcal{S}\) by Lemma 92.3, contradicting the maximality of \(N\). Therefore, \(N = M\). \(\square\)

**Corollary 92.7.** Let \(R\) be a Noetherian ring and \(M\) a nonzero finitely generated \(R\)-module. Then \(\text{Ass}_R(M)\) is a finite set.

**Proof.** By Proposition 92.6, there exists a chain of \(R\)-modules

\[
0 = M_0 < M_1 < \cdots < M_n = M
\]

satisfying \(M_i/M_{i-1} \cong R/\mathfrak{p}_i\) with \(\mathfrak{p}_i\) a prime ideal and \(\text{Ass}_R(M_i/M_{i-1}) = \{\mathfrak{p}_i\}\) for \(i = 1, \ldots, n\). By Lemma 92.3, we have \(\text{Ass}_R(M_i) \subset \text{Ass}_R(M_{i-1}) \cup \text{Ass}_R(M_i/M_{i-1})\), it follows that

\[
\text{Ass}_R(M) \subset \text{Ass}_R(M_1) \cup \text{Ass}_R(M_2/M_1) \cup \cdots \cup \text{Ass}_R(M_n/M_{n-1})
\]

(*)

\[
= \{\mathfrak{p}_1, \ldots, \mathfrak{p}_n\}.
\]

\(\square\)

**Remark 92.8.** In (*) in the above proof, it is possible that \(\text{Ass}_R M \subset \{\mathfrak{p}_1, \ldots, \mathfrak{p}_n\}\).
For example, let $M = R$ be a Noetherian domain. Then $\text{Ass}_R(R) = \{0\}$. Let $p$ be a nonzero prime ideal in $R$. Then $0 < p < R$ is a sequence of $R$-modules. If $p$ is a principal prime ideal in $R$, then $p$ is $R$-free. In particular, $p/(0) = p \cong R$ is $R$-free, giving such an example, e.g., if $R$ is a UFD (not a field).

Kaplansky believed the following corollary is one of the most useful facts.

**Corollary 92.9.** Let $R$ be a Noetherian ring, $M$ a nonzero finitely generated $R$-module, and $S$ a subrng of $R$ in $zd(M)$. Then there exists an associated prime $p$ of $M$ satisfying $S \subset p$. In particular, there exists a nonzero element $m$ in $M$ such that $Sm = 0$.

**Proof.** We have $S \subset zd(M) = \bigcup_{\text{Ass}_R(M)} p$ and $|\text{Ass}_R(M)| < \infty$, so there exists a prime ideal $p \in \text{Ass}_R(M)$ with $S \subset p$ by the Prime Avoidance Lemma 90.16. If $p = \text{ann}_R m$, then $Sm = 0$. □

**Definition 92.10.** Let $M$ be a nonzero $R$-module and $p$ an associated prime of $M$. We say that $p$ is a *minimal prime* or *isolated prime* if $p$ is minimal in $\text{Ass}_R(M)$ and an *embedded prime* otherwise.

We shall characterize minimal associated primes of a finitely generated module over a Noetherian ring below. Embedded associated primes are more mysterious. We identify one type of embedded associated primes (when they exist) in the following:

**Corollary 92.11.** Let $R$ be a Noetherian ring and $M$ a nonzero finitely generated $R$-module. Then the following sets are identical:

1. $\{p \mid p \in \text{Ass}_R(M) \text{ is maximal}\}$.
2. $\{p \in \text{Spec}(R) \mid p \subset zd(M) \text{ is maximal}\}$.
3. $\{A < R \text{ an ideal} \mid A \subset zd(M) \text{ is maximal}\}$.

**Proof.** If $A \subset zd(M)$ is a ideal, then there exists a prime ideal $p \in \text{Ass}_R(M)$ satisfying $A \subset p$ by Lemma 92.9. The result now follows easily. □

We need three lemmas to characterize minimal elements in $\text{Ass}_R(M)$ when $M$ is a finitely generated module over a Noetherian ring. We leave a proof of the first lemma as an exercise.

**Lemma 92.12.** Let $M$ be a finitely generated $R$ module and $S$ a multiplicative set in $R$. Then $\text{ann}_{S^{-1}R}(S^{-1}M) = S^{-1}(\text{ann}_R(M))$.

**Lemma 92.13.** Let $M$ be a finitely generated $R$ module over a Noetherian ring $R$ and $S$ a multiplicative set in $R$. Then

$\text{Ass}_{S^{-1}R}(S^{-1}M) = \{S^{-1}p \mid p \in \text{Ass}_R M \text{ with } p \cap S = \emptyset\}$. 
Proof. \((\supset)\): We know that
\[
\text{Spec} \left( S^{-1}(R) \right) = \{ S^{-1}p \mid p \in \text{Spec}_R(M), p \cap S = \emptyset \}
\]
and that \(\text{ann}_{S^{-1}R} \left( \frac{m}{1} \right) = S^{-1}(\text{ann}_R m)\) for all \(m\) in \(M\) by Lemma 92.12, so \(\supset\) follows.

\((\subset)\): Let \(\mathfrak{p} \in \text{Ass}_{S^{-1}R}(S^{-1}M)\), say \(\mathfrak{p} = \text{ann}_{S^{-1}R} \left( \frac{m}{1} \right)\) for some \(m \in M\) and \(s \in S\). Let \(p \in \text{Spec}(R)\) satisfy \(S^{-1}p = \mathfrak{p}\) with \(p \cap S = \emptyset\).

Claim. \(p \in \text{Ass}_R(M)\) (hence we are done).

Since \(R\) is Noetherian, \(p\) is finitely generated, say \(p = (a_1, \ldots, a_n)\). By assumption,
\[
\frac{a_i m}{1} = 0 \quad \text{in} \quad S^{-1}M,
\]
so there exist \(s_i \in S\) with \(s_i a_i m_i = 0\) for \(i = 1, \ldots, n\). Let \(\tilde{s} = s_1 \cdots s_n\), so \(\tilde{s} a_i m = 0 = a_i \tilde{s} m\) for \(i = 1, \ldots, n\). Therefore, \(p = (a_1, \ldots, a_n) \subset \text{ann}_R(\tilde{s} m)\). Let \(\varphi : R \to S^{-1}R\) be the natural ring homomorphism given by \(r \mapsto \frac{r}{1}\). Since \(\frac{s}{1}\) and \(\frac{\tilde{s}}{1}\) lie in \((S^{-1}R)^\times\), we have \(\text{ann}_{S^{-1}R} \left( \frac{\tilde{s} m}{1} \right) = \text{ann}_{S^{-1}R} \left( \frac{m}{s} \right)\).

Using Lemma 92.12, we see that
\[
\varphi^{-1}(S^{-1}(\text{ann}_R \tilde{s} m)) = \varphi^{-1}( \text{ann}_{S^{-1}R} \left( \frac{\tilde{s} m}{1} \right) ) = \varphi^{-1}( \text{ann}_{S^{-1}R} \left( \frac{m}{s} \right) ) = \varphi^{-1}(\mathfrak{p}) = p.
\]

Therefore, \(\text{ann}_R \tilde{s} m \subset \varphi^{-1}(S^{-1}(\text{ann}_R \tilde{s} m)) \subset p\). Consequently, \(p = \text{ann}_R \tilde{s} m\).

Definition 92.14. Let \(M\) be an \(R\)-module. Then the support of \(M\) is defined to be the set \(\text{Supp}(M) := \{ p \in \text{Spec} R \mid M_p \neq 0 \}\).

Lemma 92.15. Let \(R\) be a commutative ring and \(M\) and \(R\)-module.

Then
\[(1) M = 0 \text{ if and only if } \text{Supp}(M) = \emptyset.\]

\[(2) \text{If } 0 \to M' \to M \to M'' \to 0 \text{ is an exact sequence of } R\text{-modules, then } \text{Supp}_R(M) = \text{Supp}_R(M') \cup \text{Supp}_R(M'').\]

\[(3) \text{If } M \text{ is finitely generated, then } \text{Supp}(M) = V(\text{ann}_R(M)).\]

Proof. \((1)\) and \((2)\) follow by Exercise 89.31(7).

\((3)\): Let \(p \in V(\text{ann}_R(M))\) and \(S = R \setminus p\). Suppose that \(p \notin \text{Supp}(M)\), then \(M_p = 0\); so for all \(m \in M\) there exists an element \(s_m \in S\) satisfying \(s_m m = 0\). Since \(M\) is finitely generated, there exists an
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\[ s \in S \text{ satisfying } sM = 0, \text{ hence } s \in \text{ann}_R(M) \subset p, \text{ a contradiction.} \]
Therefore, \( V(\text{ann}_R(M)) \subset \text{Supp}(M) \)

Suppose that \( p \in \text{Supp}(M) \) and \( S = R \setminus p \). Then there exists an \( m \in M \) satisfying \( \frac{m}{s} \) is nonzero in \( M_p \) for some \( s \in S \). It follows that \( s \frac{m}{s} \neq 0 \) in \( M \) for all \( s \in S \). Therefore, \( S \cap \text{ann}_R(m) = \emptyset \), i.e., \( \text{ann}_R(m) \subset p \). \( \Box \)

**Theorem 92.16.** Let \( R \) be a commutative ring and \( M \) a nonzero \( R \)-module. Then \( \text{Ass}_R(M) \subset \text{Supp}(M) \). If, in addition, \( R \) is Noetherian and \( M \) is a finitely generated \( R \)-module, then the following sets are identical:

1. \( \{ p \mid p \in \text{Ass}_R(M) \text{ is minimal}\} \)
2. \( \{ p \mid p \in \text{Supp}(M) \text{ is minimal}\} \)
3. \( \{ p \mid p \in V(\text{ann}_R(M)) \text{ is minimal}\} \)

**Proof.** Let \( p \in \text{Ass}_R(M) \), say \( p = \text{ann}_R m \) for some nonzero \( m \) in \( M \). Then \( Rm \subset M \) and, as \( Rm \) is finitely generated, by Lemma 92.15,

\[ p \in V(\text{ann}_R(Rm)) = \text{Supp}(Rm) \subset \text{Supp}(M). \]

So \( \text{Ass}_R(M) \subset \text{Supp}(M) \) as needed.

Now assume that \( R \) is Noetherian and \( M \) is finitely generated. Then by Lemma 92.15, the set of elements in (2) and (3) are equal. We must show that the set of elements in (1) and (2) are equal. Since \( \text{Ass}_R(M) \subset \text{Supp}(M) \), it suffices to show:

**Claim.** If \( p \in \text{Supp}(M) \) is minimal, then \( p \in \text{Ass}_R(M) \).

By Lemma 92.13, it suffices to show that \( S^{-1}p \in \text{Ass}_{S^{-1}R}M_p \). As \( p \in \text{Supp}(M) \) is minimal, clearly \( S^{-1}p \in \text{Supp}_{S^{-1}R}(M_p) \) is also minimal. Therefore, we may assume that \( R \) is a local ring with maximal ideal of \( p \). Then \( M = M_p \) is nonzero and \( p \in \text{Supp}(M) \) is minimal, we must have \( \text{Supp}(M) = \{ p \} \). Since \( R \) is Noetherian, we have \( \emptyset \neq \text{Ass}_R(M) \subset \text{Supp}(M) = \{ p \} \) by Proposition 92.4, so \( \text{Ass}_R((M) = \{ p \} \) as needed. \( \Box \)

**Corollary 92.17.** Let \( R \) be a Noetherian ring, \( \mathfrak{a} < \mathfrak{r} \) and ideal. Then \( p \in \text{Ass}_R(R/\mathfrak{a}) \) is minimal if and only if \( p \in V(\mathfrak{a}) \) is minimal. In particular, \( \text{Min}(R) \subset \text{Ass}_R(R) \) and are the minimal elements of \( \text{Ass}_R(R) \).

**Proof.** This is immediate as \( \text{ann}_R(R/\mathfrak{a}) = \mathfrak{a} \) \( \Box \)

Note that this gives another proof that \( \text{Min}(R) \) is finite if \( R \) is Noetherian. (Cf. Corollary 89.27.)

**Corollary 92.18.** Let \( R \) be a Noetherian ring and \( M \) a finitely generated \( R \)-module. Then \( \text{Supp}(M) = \bigcup_{\text{Ass}_R(M)} V(p) \).
Proof. Minimal elements in Supp(M) and Ass_R(M) coincide. □

Remark 92.19. The Primary Decomposition Theorem when M = R is a Noetherian ring and N = \mathfrak{A} < R an ideal can now be generalized as follows — we leave the verification of the details to the reader: Let R be a commutative ring, M a finitely generated R-module, and N < M a submodule. If p \in \text{Spec}(R), we say N is p-primary if Ass_R(M/N) = \{p\}. If p \in \text{Spec}(R) and M is an R-module that is the intersection of p-primary submodules of M, then M is also p-primary.

A decomposition of N as an intersection

\begin{align*}
N = Q_1 \cap \cdots \cap Q_n
\end{align*}

with each Q_i \subset M a p_i-primary submodule, i = 1, \ldots, n, is called a primary decomposition of N in M. Such a primary decomposition is called irredundant if no Q_i can be omitted and p_i \neq p_j for i, j = 1, \ldots, n.

The Primary Decomposition Theorem generalize to: If R is a Noetherian ring, M a finitely generated R-module and N < M a submodule, then N has an irredundant primary decomposition in M. Moreover, Ass_R(M/N) = \{p_1, \ldots, p_n\}, with the p_i-primary submodule Q_i unique if p_i \in Ass_R(M/N) is minimal, i.e., p_i \in (V(\text{ann}_R(M/N))) is minimal.

Exercises 92.20.
1. Show if \mathfrak{A} is an ideal in R, then Ass_R(V(\mathfrak{A})) = Ass_R(R/\mathfrak{A}).
2. Prove Lemma 92.12.
3. Prove the assertions in Remark 92.19.
4. Let R be a Noetherian domain and M be a finitely generated free R-module. Show that Ass_R(M) = \{0\} but if M is R-torsion-free but not R-free, then the set of primes ideals appearing in the filtration of M Proposition 92.8 for M is Ass_R(M).

93. Akizuki and Krull-Akizuki Theorems

As is true in this chapter, all rings are commutative, although straight-forward modifications of the beginning of this chapter can be shown to be true if R is arbitrary.

Recall that a commutative ring is called Artinian ring if it satisfies the descending chain condition. Similarly, a left R-module is called Artinian if submodules of it satisfies the descending chain condition. We wish first with to prove a theorem of Akizuki that characterizes Artinian commutative rings as those that are Noetherian of dimension zero. We have previously noted that the analogue of the Jordan-Hölder Theorem for vector spaces holds and used it to show that dimension
was a well defined invariant for finite dimensional vector spaces. We expand on these observations.

**Definition 93.1.** Let $R$ be a commutative ring. An $R$-module $M$ is said to have a composition series if the exists a finite filtration of submodules of $M$

\[(*) \quad M = M_0 > M_1 > \cdots > M_n = 0\]

with each quotient $M_i/M_{i+1}$ irreducible (or simple), i.e., has no proper submodules. We say that a module having a composition series has finite length. The composition series in $(*)$ is called a series of length $n$, i.e., the number of links in $(*)$.

Just as in the group theoretical case, we have:

**Theorem 93.2.** (Jordan-Hölder Theorem) Suppose that an $R$-module $M$ has a composition of length $n$. Then every composition series of $M$ is of length $n$. Moreover, if $M$ has a composition series, any proper chain (i.e., all links proper) of $M$ can be refined to a composition series with unique quotients up to order an isomorphism.

We omit the proof which is similar, but easier than the group theoretic case.

If $M$ is an $R$-module having a composition series, the common length of a composition series for $M$ is called the length of $R$ and denoted by $l(M)$. We set the length of the zero module as zero. If a module $M$ has no composition series, for convenience, we let $l(M) = \infty$ (with $\infty + n = \infty$ for any integer $n$). An immediate consequence is:

**Corollary 93.3.** Let $M$ be an $R$-module having a composition series and $N$ a submodule of $M$, Then $l(N) \leq l(M)$ with equality if and only if $N = M$.

More generally,

**Corollary 93.4.** Let

\[0 \to M' \to M \to M'' \to 0\]

be an exact sequence of $R$-modules. Then $l$ is additive, i.e.,

\[l(M) = l(M') + l(M''),\]

whose proof we leave as an exercise. For a module to have finite length is a very strong condition. Indeed:

**Lemma 93.5.** Let $M$ be an $R$-module. Then $M$ has finite length if and only if it is both Noetherian and Artinian.
Proof. If $M$ has finite length, all proper chains of submodules of $M$ have length bounded by the length of $M$ by the Jordan-Hölder Theorem. Conversely, suppose that $M$ is both Noetherian and Artinian. Since it is Noetherian, there exists a maximal proper submodule $M_1 < M$, i.e., $M/M_1$ is irreducible by the Correspondence Principle. Since $M$ is Noetherian so is $M_1$. Continuing gives a descending chain of modules $M = M_0 > M_1 > M_2 > \cdots$ with $M_i/M_{i+1}$ irreducible. This must stop, since $M$ is also Artinian. □

An easy consequence of this (that is left as an exercise) is:

**Proposition 93.6.** Let $V$ be a vector space over a field $K$. Then the following are equivalent:

1. $V$ is finite dimensional.
2. $V$ has finite length.
3. $V$ is Noetherian.
4. $V$ is Artinian.

Moreover, the dimension of $V$ is just its length.

**Corollary 93.7.** Suppose $m_1, \ldots, m_n$ are maximal ideals in $R$ (not necessarily distinct) and further that $m_1 \cdots m_n = 0$. Then $R$ is Noetherian if and only if $R$ is Artinian.

**Proof.** Each $m_1 \cdots m_i/m_1 \cdots m_{i-1}$ is an $(R/m_i)$-vector space so is Artinian if and only if it is Noetherian. We know that $m_1 \cdots m_i$ is Noetherian (resp., Artinian) if and only if both of the ideals $m_1 \cdots m_{i+1}$ and $m_1 \cdots m_i/m_1 \cdots m_{i+1}$ are. [Cf. Proposition 39.5 and Exercise 39.12(4).] The result now follows, since

$$R \supset m_1 \supset m_1m_2 \supset \cdots \supset m_1 \cdots m_n = 0.$$  

**Proposition 93.8.** Let $R$ be a nonzero commutative Artinian ring. Then $R$ is semi-local of dimension zero. In particular, $\text{rad}(R) = \text{nil}(R)$.

**Proof.** Let $p$ be a prime ideal of $R$. Then $R/p$ is an Artinian domain. Thus to show that $\text{dim} R$ is zero, i.e., $\text{Spec}(R) = \text{Max}(R)$, we need only show that any Artinian domain is a field. But if $R$ is an Artinian domain and $x$ nonzero in $R$, then the descending chain

$$Rx \supset Rx^2 \supset Rx^3 \supset \cdots$$

must stabilize. Hence $Rx^n = Rx^{n+1}$ for some $n$. In particular, $yx^{n+1} = x^n$ for some $y$ in the domain $R$. Thus $yx = 1$ and $R$ is a field.

Now suppose that $R$ is an arbitrary Artinian ring. By the Minimal Principle, there exists a minimal element in the set of finite intersections of maximal ideals,

$$\{m_1 \cap \cdots \cap m_n \mid m_1, \ldots, m_m \in \text{Max}(R), \text{ some } m\}.$$
Let $m_1 \cap \cdots \cap m_n$ be such a minimal element. If $m$ is a maximal ideal, then by minimality

$$m_1 \cap \cdots \cap m_n = m \cap m_1 \cap \cdots \cap m_n \subset m.$$ 

It follows that $m = m_i$, for some $i$, hence $\text{Max}(R) = \{m_1, \ldots, m_n\}$. □

**Theorem 93.9.** (Akizuki) Let $R$ be a nonzero commutative ring. Then the following are equivalent.

1. $R$ is Artinian.
2. $R$ is Noetherian of dimension zero.
3. Every finitely generated $R$-module has finite length.

**Proof.** Every module over an Artinian (resp., Noetherian) ring $R$ is Artinian (resp., Noetherian), since it is a quotient of $R^n$ for some $n$. Thus by the above, we know that $R$ satisfies both (1) and (2) if and only if $R$ satisfies (3). So we need only show that $R$ satisfies (1) if and only if $R$ satisfies (2).

Suppose that (1) holds, i.e., that $R$ is Artinian. Then $R$ is semi-local of dimension zero. Let $\text{Max}(R) = \{m_1, \ldots, m_n\}$. By the corollary above, it suffices to show that $0 = \prod_{i=1}^n m_i^k$ for some $k$. But $\prod_{i=1}^n m_i \subset \bigcap_{i=1}^n m_i = \text{rad}(R) = \text{nil}(R)$ by the proposition. So it suffices to prove the following:

**Claim 93.10.** If $R$ is Artinian, then $\text{nil}(R)$ is nilpotent:

By the descending chain condition, $(\text{nil}(R))^k = (\text{nil}(R))^{k+i}$ for some $k$ and all $i$. Let $\mathfrak{A} := (\text{nil}(R))^k$. We must show that $\mathfrak{A} = 0$. Suppose not. Let

$$S = \{\mathfrak{B} < R | \mathfrak{B} \text{ an ideal of } R \text{ with } \mathfrak{A}\mathfrak{B} \neq 0\}.$$ 

Since $\mathfrak{A} \in S$ and $R$ is Artinian, there exists an ideal $\mathfrak{B}$ that is minimal in $S$. In particular, there exists an $x$ in $\mathfrak{B}$ such that $x\mathfrak{A} \neq 0$. By minimality, $\mathfrak{B} = Rx$. Since $(x\mathfrak{A})^2 = x^2\mathfrak{A} = x\mathfrak{A} \neq 0$, we also have $x\mathfrak{A} = \mathfrak{B}$ by minimality. Choose $y \in \mathfrak{A}$ such that $x = xy$. Then $x = xy^N$ for all positive integers $N$. But $y \in \mathfrak{A} \subset \text{nil}(R)$, so $y^N = 0$ for some $N$ and hence $x = 0$ also. This is a contradiction. Thus $\mathfrak{A} = 0$ and the Claim is established.

Now suppose that $R$ is Noetherian of dimension zero. By the corollary above, it suffices to show that there exist maximal ideals $m_1, \ldots, m_n$ in $R$, not necessarily distinct, so that $0 = m_1 \cdots m_n$. Since $R$ is Noetherian the zero ideal contains a finite product of prime ideals. (Cf. Exercise 29.22(14).) Since $\dim R = 0$, these prime ideals are maximal. The result follows. □
Corollary 93.11. Let $R$ be a domain. Then the following are equivalent:

(1) $R$ is Noetherian of dimension at most one.
(2) If $0 < \mathfrak{A} < R$ is an ideal then $R/\mathfrak{A}$ has finite length.
(3) If $0 < \mathfrak{A} < R$ is an ideal then $R/\mathfrak{A}$ is Artinian.

Proof. If (1) holds and $0 < \mathfrak{A} < R$ is an ideal then $R/\mathfrak{A}$ is Noetherian of dimension zero so (2) holds. Clearly, (2) implies (3), so we need only show that (3) implies (1).

If (3) holds then $R/\mathfrak{A}$ is a Noetherian ring of dimension zero for any ideal $0 < \mathfrak{A} < R$. In particular, if $x \in \mathfrak{A}$ is nonzero, then $\mathfrak{A}/Rx$ is a finitely generated ideal in $R/Rx$. It follows that $\mathfrak{A}$ is finitely generated as an ideal in $R$, i.e., $R$ is Noetherian. If $0 < p_1 \subset p_2$ is a chain of primes in $R$ then $p_2/p_1 = 0$ in the Artinian domain, hence field, $R/p_1$. Thus $R$ has dimension at most one.

Lemma 93.12. Let $M$ be a non-trivial $R$-module and $p$ a prime ideal in $R$ containing $\operatorname{ann}_R(M)$ and a minimal such prime ideal, i.e., no prime ideal containing $\operatorname{ann}_R(M)$ properly lies in $p$. Then $p$ consists of zero divisors of $M$. In particular,

$$\bigcup_{\operatorname{Min}(R)} p \subset \operatorname{zd}(R).$$

[There is no Noetherian condition or finite generation condition.]

Proof. Let $S$ be the multiplicative set in $R$ defined by

$$S := \{ab \mid a \in R \setminus p \text{ and } b \in R \setminus \operatorname{zd}(M)\}.$$

Claim 93.13. $S \cap \operatorname{ann}_R(M) = \emptyset$.

Suppose not. Then there exists an $a$ in $R \setminus p$ and $b$ in $R \setminus \operatorname{zd}(M)$ satisfying $abM = 0$. Since $b$ is not a zero divisor on $M$, we have $aM = 0$, and hence $a \in \operatorname{ann}_R(M) \subset p$, a contradiction. This establishes the Claim.

Thus there exists a prime $\mathfrak{P}$ containing $\operatorname{ann}_R(M)$ such that $\mathfrak{P}$ excludes $S$ and is maximal with respect to this property. Since 1 lies in $R$ but not in $p$ or $\operatorname{zd}(M)$ and $S = (R \setminus \operatorname{zd}(M)) \cdot (R \setminus p)$, we have

$$S \supset (R \setminus \operatorname{zd}(M)) \quad \text{and} \quad S \supset (R \setminus p).$$

Therefore,

$$\mathfrak{P} \subset \operatorname{zd}(M) \cap p \subset p.$$

The minimality condition on $p$ implies that $p = \mathfrak{P}$, so $p \subset \operatorname{zd}(M)$ as desired. For the last statement, let $M = R$. Then every prime contains
0 = \text{ann}_R(R)$. It follows from the first part that if $p$ is a minimal prime ideal in $R$, then $p$ consists of zero divisors of $R$ as $\text{ann}_R(R) = 0$.

We need the lemma in the following special case.

**Corollary 93.14.** If $\dim R = 0$ then $R^\times = R \setminus \text{zd}(R)$.

**Proof.** We have

\[(*) \quad \text{zd}(R) \supset \bigcup_{\text{Min}(R)} p = \bigcup_{\text{Spec}(R)} p = \bigcup_{\text{Max}(R)} m.\]

Since $R^\times = R \setminus \bigcup_{\text{Max}(R)} m$, in the above $(*)$ is an equality. □

**Lemma 93.15.** Let $R$ be a Noetherian domain of dimension one. Let $a$ and $c$ be non-zero elements of $R$. Let

\[\mathfrak{A} = \bigcup_{n=0}^\infty (Rc : Ra^n) := \{x \in R \mid xa^n \in Rc \text{ for some integer } n\}.\]

Then

\[\mathfrak{A} + Ra = R.\]

**Proof.** Let $\mathfrak{A}_k = (Rc : Ra^k) := \{x \in R \mid xa^k \in Rc\}$. Since $\mathfrak{A}_k \subset \mathfrak{A}_{k+1}$, for all $k$, we know that $\mathfrak{A}$ is an ideal. Since $R$ is Noetherian, there exists an integer $n$ such that $\mathfrak{A}_n = \mathfrak{A}_{n+i} = \mathfrak{A}$ for all positive integers $i$. As $c$ lies in $\mathfrak{A}_k$ for all $k$, we have $c \in \mathfrak{A}$. In particular, $\mathfrak{A}$ is not trivial. Let $\gamma : R \rightarrow R/\mathfrak{A}$ be the canonical epimorphism. Since $\mathfrak{A} > 0$ and $R$ is a domain, it is clear that $\dim R/\mathfrak{A} = 0$. [We do not need Akizuki’s Theorem.] By the corollary above, it suffices to establish the following:

**Claim 93.16.** $\overline{a}$ is not a zero divisor in $R/\mathfrak{A}$:

If this were false, then there would exist a $y$ in $R \setminus \mathfrak{A}$ satisfying $\overline{ay} = 0$, i.e., $ay \in \mathfrak{A} = \mathfrak{A}_n$. We would then have $(ay)a^n \in Rc$ and that would imply that $y \in \mathfrak{A}_{n+1} = \mathfrak{A}_n = \mathfrak{A}$, a contradiction. This establishes the Claim. □

**Theorem 93.17.** (Krull-Akizuki Theorem) Let $A$ be a Noetherian domain of dimension one. Let $F$ be the quotient field of $A$ and let $K/F$ be a finite field extension. Let $B$ be a ring such that $A \subset B \subset K$. Then $B$ is a Noetherian domain of dimension at most one and if $\mathfrak{B}$ is a non-trivial ideal of $B$, then $B/\mathfrak{B}$ is a finitely generated $(A/(\mathfrak{B} \cap A))$-module of finite length.
Proof. We first make two reductions.

**Reduction 1.** We may assume that \( F = K \):

Certainly, we may assume that \( K \) is the quotient field of \( B \) and, in fact, \( K = F(x_1, \ldots, x_n) \) for some \( x_i \in B \). Choose \( 0 \neq c \in A \) such that each \( cx_i \) is integral over \( A \). Let \( C = A[cx_1, \ldots, cx_n] \). Then \( C \) is integral over \( A \) and a finitely generated \( A \)-module. Thus \( C \) is a Noetherian domain of dimension one with quotient field \( K \) and contained in \( B \). If \( C \) is an ideal in \( C \) then \( C/C \) is integral over \( A/(C \cap A) \) and a finitely generated \( (A/(C \cap A)) \)-module. So \( C/C \) and \( A/(C \cap A) \) are Noetherian rings of the same dimension. In particular, one is Artinian if and only if the other is. Consequently, if \( B \) is a non-trivial ideal in \( B \) then by clearing denominators and multiplying by an appropriate power of \( c \), we see that \( B \cap C \) is a non-trivial ideal in \( C \). In particular, \( B/B \) is a finitely generated \( (A/(B \cap A)) \)-module if it is a finitely generated \( (C/(B \cap C)) \)-module and has finite length as an \( (A/(B \cap A)) \)-module if it has finite length as a \( (C/(B \cap C)) \)-module. This completes the reduction.

**Reduction 2.** It suffices to show that the \( (A/Aa) \)-module \( B/Ba \) is finitely generated for any \( 0 \neq a \in A \):

To show that \( B \) is Noetherian of dimension at most one, it suffices, by Corollary 93.11 to Akizuki’s Theorem, to show that \( B/B \) is Artinian for any nonzero ideal \( B \) of \( B \). Let \( 0 < B \) be an ideal of \( B \). We must also show that \( B/B \) has finite length over \( A/(B \cap A) \).

**Claim.** \( 0 < B \cap A \):

Let \( B' \) be any non-trivial finitely generated \( A \)-submodule of \( B \). Then there exists \( 0 \neq c \in A \) such that \( cB' \) lies in the domain \( A \) by the first reduction. This establishes the Claim.

Let \( 0 \neq a \) lie in \( B \cap A \). Then \( A/Aa \) is Artinian by Corollary 93.11 to Akizuki’s Theorem. By assumption, \( B/Ba \) is a finitely generated \( (A/Aa) \)-module. Since \( B/B \) is a cyclic \( (B/Ba) \)-module, it is also finitely generated as an \( (A/Aa) \)-module, hence has finite length over the Artinian ring \( A/Aa \) so also over the Artinian ring \( A/(B \cap A) \). This also implies that \( B/B \) is Artinian.

So to finish we are in the following situation. We have \( 0 \neq a \in A \) is fixed, and we must show that \( B/Ba \) is finitely generated as an \( (A/Aa) \)-module. We do this in a number of steps. Note, as before, that \( A/Aa \) is an Artinian ring.

**Step 1.** Let \( x \in B (\subseteq F) \). Then there exists a positive integer \( n \) such that \( x \in Aa^{-n} + Ba \):
Write \( x = \frac{b}{c} \) with \( b, c \in A \) and \( c \neq 0 \). Set

\[
\mathfrak{B} = \bigcup_{n=0}^{\infty} (Ac : Aa^n) := \{ y \in A \mid ya^n \in Ac, \text{ some } n \}.
\]

By Lemma 93.15, we have \( \mathfrak{B} + Aa = A \) so \( 1 = y + za \), some \( y \in \mathfrak{B} \) and \( z \in A \). Consequently, \( x = yx + zax \). Since \( y \in \mathfrak{B} \), by definition, there exists an integer \( n \) such that \( ya^n \in Ac \), so

\[
x = yx + zax = \frac{ya^n b}{a^n c} + zax \text{ lies in } Aa^{-n} + Ba
\]
as needed.

**Step 2.** Let \( \mathfrak{A}_n := (Ba^n \cap A) + Aa \), an ideal of \( A \). Then there exists a positive integer \( m \) such that \( A^m = A^m + i \) for all positive integers \( i \):

Each \( \mathfrak{A}_n \) contains \( a \) so is non-trivial. Moreover, it is clear that the \( \mathfrak{A}_n \) form a descending chain of ideals. Since \( A/Aa \) is Artinian, the descending chain of ideals

\[
\cdots \supset \mathfrak{A}_n/Aa \supset \mathfrak{A}_{n+1}/Aa \supset \cdots
\]
stabilizes, hence so does the corresponding chain of \( \mathfrak{A}_n \)'s.

**Step 3.** Let \( m \) be the integer in Step 2. Then \( B \subset Aa^{-m} + Ba \):

Let \( x \in B \) be fixed. Then by Step 1, there exists a minimal positive integer \( n \) so that \( x \in Aa^{-n} + Ba \). If we show that \( m \geq n \) then \( a^m \in Aa^n \) and \( Aa^{-n} \subset Aa^{-m} \) as needed. So we may assume that \( n > m \). Write \( x = ra^{-n} + ba \), with \( r \in A \) and \( b \in B \). Then

\[
r = a^n(x - ba) \in Ba^n \cap A \subset (Ba^n \cap A) + Aa = \mathfrak{A}_n
\]
and \( \mathfrak{A}_n = \mathfrak{A}_{n+1} = \mathfrak{A}_m \). Hence \( r = b_1a^{n+1} + r_1a \), for some \( r_1 \in A \) and \( b_1 \in B \) so \( x = ra^{-n} + ba = (b_1a^{n+1} + r_1a)a^{-n} + ba \) lies in \( Aa^{-n+1} + Ba \).

This contradicts the minimality of \( n \), so completes the step.

**Step 4.** \( B/Ba \) is a finitely generated \((A/Aa)\)-module (and hence we are done):

By Step 3, we know that \( B/Ba \subset (Aa^{-m} + Ba)/Ba \). Moreover, we know that the \( A \)-module \((Aa^{-m} + Ba)/Ba \cong Aa^{-m}/(Aa^{-m} \cap Ba)\) is cyclic, hence Noetherian. Thus \( B/Ba \) is a finitely generated \( A \)-module as needed.

**Corollary 93.18.** Let \( A \) be a Dedekind domain with quotient field \( F \).
If \( K/F \) is a finite field extension then \( A_K \) is also a Dedekind domain.

**Remarks 93.19.** 1. If \( K = F \) in the theorem then clearly \( F \) is the only field between \( A \) and \( F \), i.e., the only such \( B \) with \( \dim B = 0 \) is \( F \).
2. If \( B \) is zero and \( B = F \) then \( F \) is not a finitely generated \( A \)-module when \( A < F \).

**Exercises 93.20.**

1. Prove Corollary 93.4
2. Prove Proposition 93.6
3. Let \( R \) be a commutative Artinian ring and \( M \) a finitely generated free \( R \)-module. Let \( N \) be a submodule of \( M \) that is also \( R \)-free. Show that \( \text{rank} N \leq \text{rank} M \).
4. Let \( R \) be a commutative ring and \( M \) a finitely generated free \( R \)-module. Let \( N \) be a submodule of \( M \) that is also \( R \)-free. Show that \( \text{rank} N \leq \text{rank} M \).
5. Why does it suffice to prove Claim 93.16.
6. Let \( A \) be a domain with quotient field \( F \). Suppose that every ring \( B \) with \( A \subset B \subset F \) is Noetherian. Show that \( \dim A = 1 \).

### 94. Affine Algebras

Throughout this section \( F \) will denote a field and \( R \) a commutative ring.

Recall a finitely generated commutative \( F \)-algebra, with \( F \) a field, is called an **affine \( F \)-algebra**. By the Hilbert Basis Theorem, such an algebra is a Noetherian ring. We wish to investigate the dimension theory attached to such algebras. We know that an integral extension of rings preserves dimension. As an affine \( F \)-algebra \( A \) is isomorphic to a quotient of a polynomial ring \( F[t_1, \ldots, t_n] \) for some \( n \), it seems reasonable to assume that if \( A \) is also a domain, then the transcendence degree of \( qf(A)/F \) should play a crucial role in such an investigation.

A fundamental theorem (whose proof we give is due to Nagata) about affine \( F \)-algebras bringing together these two observations is the following:

**Theorem 94.1.** (Noether Normalization Theorem) Let \( F \) be a field and \( A = F[x_1, \ldots, x_n] \) an affine \( F \)-algebra that is also a domain. Let \( r = \text{tr} \deg_F qfA \). Then there exist \( y_1, \ldots, y_r \) in \( A \) algebraically independent over \( F \) satisfying \( A/F[y_1, \ldots, y_r] \) is integral. In particular, \( A \) is a finitely generated \( F[y_1, \ldots, y_r] \)-module. Let \( \varphi : F[t_1, \ldots, t_r] \to A \) be the evaluation map given by \( t_i \mapsto y_i \). Then

1. \( \varphi \) is an \( F \)-algebra monomorphism.
2. \( ^*\varphi : \text{Spec}(A) \to \text{Spec}(F[t_1, \ldots, t_r]) \) is surjective and has finite fibers.
3. \( \dim A = \dim F[y_1, \ldots, y_r] = \dim F[t_1, \ldots, t_r] \).
Proof. Let \( A = F[t_1, \ldots, t_n] \). If \( qf(A)/F \) is a finite field extension, then it is integral. Consequently, we may assume that \( qf(A)/F \) is not algebraic. Relabeling, we may also assume that \( x_1 \) is transcendental over \( F \). In addition, we may assume that \( x_1, \ldots, x_n \) are algebraically dependent over \( F \).

Let \( (j) := (j_1, \ldots, j_n) \) with \( j_i \) non-negative integers for \( i = 1, \ldots, n \). As the \( x_i \) are algebraically dependent, we have an equation

\[
(*) \quad \sum_{(j)} a_{(j)} x_1^{j_1} \cdots x_n^{j_n} = 0
\]

for some \( a_{(j)} = a_{j_1, \ldots, j_n} \in F \) not all zero. Let \( m_2, \ldots, m_n \) be positive integers (to be chosen) and set

\[
y_i = x_i - x_1^{m_i} \quad \text{for} \quad i = 2, \ldots, n
\]

and \( m_1 = 1 \). Define

\[
(j) \cdot (m) := j_1 + j_2 m_2 + \cdots + j_n m_n.
\]

Plugging \( x_i = y_i + x_1^{m_i}, \ i = 2, \ldots, n \), into \((*)\) yields an equation

\[
(**) \quad \sum a_{(j)} x_1^{(j) \cdot (m)} + f(x_1, y_2, \ldots, y_n) = 0
\]

with \( f \in F[x_1, y_2, \ldots, y_n] \) and having no monomial solely in \( x_1 \). In addition, if there is no cancellation in the term \( \sum a_{(j)} x_1^{(j) \cdot (m)} \), then then degree in \( x_1 \) of the equation in \((**)\) satisfies.

\[
(\dagger) \quad \deg_{x_1} f < \deg_{x_1} \left( \sum a_{(j)} x_1^{(j) \cdot (m)} \right).
\]

Choose \( d > \max\{j_i \mid a_{(j)} \neq 0\} \), and set

\[
(m) = (1, d, \ldots, d^{m-1}), \ i.e., \ m_i = d^{i-1} \quad \text{for} \quad i = 0, \ldots, n - 1.
\]

As the \((j)\)'s are distinct \( n \)-tuples, the \((j) \cdot (m)\)'s are all distinct, since they are the \( d \)-adic expansions of the \((j) \cdot (m)\) and \( x_1 \) is transcendental over \( F \). Thus \((\dagger)\) holds, so \((**)\) produces an integral equation for \( x_1 \) over \( F[y_2, \ldots, y_n] \), as the leading coefficient of \((***)\) in \( x_1 \) say \( a_{(j)} \) is a unit in \( F \). Since \( x_i - (x_1^{m_i} + y_i) = 0 \) for \( i = 2, \ldots, n \) by definition, \( x_i \) is integral over \( F[x_1, y_2, \ldots, y_n] \) for \( i > 1 \), hence over \( F[y_2, \ldots, y_n] \). Consequently, \( F[x_1, x_2, \ldots, x_n]/F[y_2, \ldots, y_n] \) is integral and \( F[x_1, x_2, \ldots, x_n] \) is a finitely generated \( F[y_2, \ldots, y_n] \)-module.

If \( y_2, \ldots, y_n \) are algebraically independent over \( F \), we are done. If not we can repeat the process to finish this part of the theorem (as being integral and being a finitely generated module persists in towers).
The other statements in the theorem follow from the above and our previous work, e.g., using Lying Over, as \( qf(A/F(y_1, \ldots, y_r)) \) is a finite algebraic extension and \( F[y_1, \ldots, y_r] \cong F[t_1, \ldots, t_r] \) is a normal domain. □

**Remark 94.2.** Using the Hilbert Nullstellensatz (cf. the discussion following Lemma 40.14 or another proof of it in 94.10 below), a geometric picture for the Noether Normalization Theorem for a curve when \( F = \mathbb{C} \) is given by the following picture, where \( f(x, y) = 0 \) is the curve in \( \mathbb{C}^2 \) given by \( f = y^2 - x^2(x + 1) \) in \( \mathbb{C}[x, y] \) with \( x, y \) variables the nodal cubic and \( A \) is the affine \( \mathbb{C} \)-algebra \( \mathbb{C}[x, y]/(f) \):

\[
\begin{array}{c}
\text{Max}(A) \\
A = \mathbb{C}[x, y]/(f) \\
\text{Spec}(\mathbb{C}[t]) \setminus \{(0)\} \\
\mathbb{C}[t]
\end{array}
\]

**Corollary 94.3.** Let \( F \) be a field. Then
\[
\dim F[t_1, \ldots, t_n] = n = \text{tr deg}_F F(t_1, \ldots, t_n).
\]

**Proof.** Since \( (0) < (t_1) < \cdots < (t_1, \ldots, t_n) \) is a chain of prime ideals in \( F[t_1, \ldots, t_n] \), we have \( \dim F[t_1, \ldots, t_n] \geq n \). Suppose that \( (0) < \mathfrak{p}_1 \subset \cdots \subset \mathfrak{p}_m \) is a chain of prime ideals in \( F[t_1, \ldots, t_n] \) and \( A = F[t_1, \ldots, t_n]/\mathfrak{p}_1 \), an affine \( F \)-domain, with \( - : F[t_1, \ldots, t_n] \to A \) the canonical epimorphism, so \( A = F[t_1, \ldots, t_m] \). By Noether Normalization, there exist \( y_1, \ldots, y_r \) in \( A \), algebraically independent over \( F \) satisfying \( A/F[y_1, \ldots, y_r] \) is integral. Since \( 0 < \mathfrak{p}_1 \), we must have \( \overline{t_1}, \ldots, \overline{t_n} \) are algebraically dependent over \( F \), i.e., \( r < n \). By induction, we have
\[
\dim A = \dim F[y_1, \ldots, y_r] = \dim F[t_1, \ldots, t_r] = r = \dim F[t_1, \ldots, t_n]/\mathfrak{p}_1 \geq m - 1,
\]

since \( \overline{\mathfrak{p}_1} \subset \cdots \subset \overline{\mathfrak{p}_m} \) is a chain of prime ideals in \( A \). Consequently, \( n \geq m \), and the result follows. □

**Notation 94.4.** Let \( A \) be an affine \( F \)-algebra that is also a domain and \( K = qf(A) \). Then there exists a transcendence basis of \( K \) over \( F \)
in $A$ by clearing denominators (or using Noether Normalization). We shall set
$$\text{tr} \deg_F A := \text{tr} \deg_F K.$$ 

**Corollary 94.5.** Let $A$ be an affine $F$-algebra that is also a domain. Then $\dim A = \text{tr} \deg_F A$.

**Proof.** Let $r' = \text{tr} \deg_F A$. By Noether Normalization there exist $y_1, \ldots, y_r$ in $A$ algebraically independent over $F$ with $A/F[y_1, \ldots, y_r]$ integral. Therefore, $qf(A)/F(y_1, \ldots, y_r)$ is finite algebraic, so $r = r'$, and we have
$$\dim A = \dim F[y_1, \ldots, y_r] = \text{tr} \deg_F F[y_1, \ldots, y_r] = r = \text{tr} \deg_F A.$$ 

**Corollary 94.6.** Let $A$ be an affine $F$-algebra and $\mathfrak{p}$ a prime ideal in $A$. Then $\dim V(\mathfrak{p}) = \text{tr} \deg_F A/\mathfrak{p}$.

**Proof.** $\dim V(\mathfrak{p}) = \dim A/\mathfrak{p}$. 

As a consequence of Noether Normalization, we also obtain another proof of Zariski’s Lemma.

**Corollary 94.7.** (Zariski’s Lemma) Let $A$ be an affine $F$-algebra. If $A$ is a field, then $A/F$ is a finite field extension.

**Proof.** We have $0 = \dim A = \text{tr} \deg_F A$, so the finitely generated field extension $A/F$ is algebraic, hence finite. 

**Corollary 94.8.** Let $\varphi : A \to B$ be an $F$-algebra homomorphism of affine $F$-algebras. Then the restriction of $^a\varphi$ to $\text{Max}(B)$ satisfies
$$^a\varphi|_{\text{Max}(B)} : \text{Max}(B) \to \text{Max}(A).$$

**Proof.** Let $\mathfrak{m}$ be a maximal ideal in $B$. Then $\mathfrak{p} = ^a\varphi(\mathfrak{m})$ is a prime ideal in $A$ and $\varphi$ induces a ring monomorphism $A/\mathfrak{p} \to B/\mathfrak{m}$. As the quotient of a affine $F$-algebra is an affine $F$-algebra, $B/\mathfrak{m}$ is a finite field extension of $F$ by by Zariski’s Lemma. Thus $F \subset A/\mathfrak{p} \subset B/\mathfrak{m}$ which implies that $A/\mathfrak{p}$ is a field. Therefore, $\mathfrak{p}$ a maximal ideal in $A$.

Recall from §35, if $A = F[t_1, \ldots, t_n]$ and $\mathfrak{A} \subset A$ is an ideal, the **affine variety** of $\mathfrak{A}$ in $F^n$ is defined by
$$Z_F(\mathfrak{A}) = \{ \mathbf{a} = (a_1, \ldots, a_n) \in F^n \mid f(\mathbf{a}) = 0 \text{ for all } f \in \mathfrak{A} \},$$
the collection of all such forms a system of closed sets for the Zariski topology of $F^n$. We have seen in Section 40, how Zariski’s Lemma implies the (Weak) Hilbert Nullstellensatz: 


Theorem 94.9. (Hilbert Nullstellensatz) (Weak Form) Suppose that $F$ is an algebraically closed field, $A = F[t_1, \ldots, t_n]$, and $\mathfrak{A} = (f_1, \ldots, f_r)$ is an ideal in $A$. Then $Z_F(\mathfrak{A})$ is the empty set if and only if $\mathfrak{A}$ is the unit ideal. In particular, if $\mathfrak{A} < R$, then there exists a point $\bar{a} \in F^n$ satisfying $f_1(\bar{a}) = 0$, $\ldots$, $f_r(\bar{a}) = 0$.

and using the Rabinowitch Trick, how this implies the (Strong) Hilbert Nullstellensatz):

Theorem 94.10. (Hilbert Nullstellensatz) (Strong Form) Suppose that $F$ is an algebraically closed field and $A = F[t_1, \ldots, t_n]$. Let $f, f_1, \ldots, f_r$ be elements in $A$ and $A = (f_1, \ldots, f_r) \subset R$. Suppose that $f(\bar{a}) = 0$ for all $\bar{a} \in Z_F(\mathfrak{A})$. Then there exists an integer $m$ such that $f^m \in \mathfrak{A}$, i.e., $f \in \sqrt{\mathfrak{A}}$. In particular, if $\mathfrak{A}$ is a prime ideal, then $f \in \mathfrak{A}$.

We give another proof of Strong Form of the Hilbert Nullstellensatz. We begin this proof by reducing it to another form of the theorem.

Proof. Let $A = F[t_1, \ldots, t_n]$ and $\operatorname{Max}(A) \subset \operatorname{Spec}(A)$ have the induced topology. The Weak Form of the Hilbert Nullstellensatz implies that the map

$$F^n \to \operatorname{Max}(A)$$

given by $\bar{a} = (a_1, \ldots, a_n) \mapsto \mathfrak{m}_\bar{a} := (t_1 - a_1, \ldots, t_n - a_n)$ is a homeomorphism. Let $\mathfrak{A} < A$ be an ideal. Recall that $\mathcal{I}(V) := \bigcap_V \mathfrak{p}$ if $V \subset \operatorname{Spec}(A)$, so $\sqrt{\mathfrak{A}} = \mathcal{I}(V(\mathfrak{A})) = \bigcap_{V(\mathfrak{A})} \mathfrak{p}$. Therefore, it suffices to prove that

$$\sqrt{\mathfrak{A}} = \bigcap_{V(\mathfrak{A})} \mathfrak{p} = \bigcap_{V(\mathfrak{A}) \cap \operatorname{Max}(A)} \mathfrak{m}$$

in the above.

In fact, this is independent of the fact that $F$ be algebraically closed and is an algebraic version of the Strong Form of the Hilbert Nullstellensatz that we shall now state and prove.

Theorem 94.11. (Hilbert Nullstellensatz) (Algebraic Form) Let $A$ be an affine $F$-algebra and $\mathfrak{A} < A$ an ideal. Then

$$\sqrt{\mathfrak{A}} = \bigcap_{V(\mathfrak{A}) \cap \operatorname{Max}(A)} \mathfrak{m},$$

i.e., the closed points in $\operatorname{Spec}(A)$ determines the variety $V(\mathfrak{A})$.

Proof. Let $\mathfrak{B} = \bigcap_{V(\mathfrak{A}) \cap \operatorname{Max}(A)} \mathfrak{m}$. Therefore, $\sqrt{\mathfrak{A}} \subset \mathfrak{B}$. Suppose that $\sqrt{\mathfrak{A}} < \mathfrak{B}$. Let $b \in \mathfrak{B} \setminus \sqrt{\mathfrak{A}}$ and $S = \{b^n \mid n \geq 0\}$. Then $S \bigcap \sqrt{\mathfrak{A}} = \emptyset$. Let $\varphi : A \to S^{-1}A$ be the canonical homomorphism given by $a \mapsto \frac{a}{1}$. Since $S^{-1}A = A[b^{-1}]$, it is also an affine $F$-algebra, hence $\varphi$ is an $F$-algebra homomorphism of affine $F$-algebras, so $a \varphi$ takes maximal ideals
to maximal ideals. By choice, \( S^{-1}\mathfrak{A} < S^{-1}A \), so there exists a maximal ideal \( n \) in \( V(S^{-1}\mathfrak{A}) \cap \text{Max}(S^{-1}A) \). Since \( \varphi(b) \) is a unit in \( S^{-1}A \), we know that \( \varphi(b) \notin n \). We also have \( \mathfrak{A} = \varphi^{-1}(S^{-1}\mathfrak{A}) \subset \varphi^{-1}(n) = a\varphi(n) \), so \( a\varphi(n) \) lies in \( V(\mathfrak{A}) \cap \text{Max}(A) \). Therefore, we have

\[
b \in \mathfrak{B} = \bigcap_{V(\mathfrak{A}) \cap \text{Max}(A)} m \subset a\varphi(n),
\]

hence \( \varphi(b) \in n \), a contradiction. \( \square \)

A ring satisfying the conclusion of this form of the Nullstellensatz is called a Jacobson (or Hilbert) ring.

**Corollary 94.12.** Let \( A \) be an affine \( F \)-algebra. Then \( \text{Max}(A) \) is a dense subset of \( \text{Spec}(A) \).

**Proof.** Let \( \mathfrak{A} < A \) be an ideal and \( U = \text{Spec}(A) \setminus V(\mathfrak{A}) \). If \( \text{Max}(A) \subset V(\mathfrak{A}) \), then we have

\[
\text{nil}(A) = \bigcap_{\text{Spec}(A)} \mathfrak{p} = \bigcap_{\text{Max}(A)} m = \bigcap_{\text{Max}(A) \cap V(\mathfrak{A})} m = \sqrt{\mathfrak{A}}
\]

by the Algebraic Form of the Hilbert Nullstellensatz, so \( \text{Spec}(A) = V(\sqrt{\mathfrak{A}}) = V(\mathfrak{A}) \) and \( U \) is empty. The result follows. \( \square \)

Let \( A \) be an affine \( F \)-algebra. Then \( A \cong F[t_1, \ldots, t_N]/\mathfrak{A} \) for some \( N \) and ideal \( \mathfrak{A} < F[t_1, \ldots, t_N] \). As \( V(\mathfrak{A}) \cong \text{Spec}(F[t_1, \ldots, t_N]/\mathfrak{A}) \cong \text{Spec}(A) \) and \( A \) is a a Noetherian ring by the Hilbert Basis Theorem, \( \text{Spec}(A) \) is a Noetherian space. Therefore, it decomposes into finitely many indecomposable components, i.e., the irreducible varieties \( V(\mathfrak{p}) \) with \( \mathfrak{p} \) in the finite set \( \text{Min}(A) \). Of course, each prime ideal (respectively, minimal prime ideal, maximal prime ideal) \( \mathfrak{p} \) in \( \text{Spec}(A) \) is isomorphic to \( \mathfrak{P}/\mathfrak{A} \) for some prime ideal \( \mathfrak{P} \) (and respectively, minimal, maximal) in \( V(\mathfrak{A}) \). If \( \mathfrak{p} \in \text{Min}(A) \), then \( \dim V(\mathfrak{p}) = \text{tr deg}_FA/\mathfrak{p} \). Different minimal primes can result in different dimensions. Of course, one is interested in determining when all irreducible components have the same dimension, a subject that we shall not pursue. Suppose that \( \mathfrak{P} \) is a prime ideal in \( A \). We define \( \text{codim}_{\text{Spec}(A)} V(\mathfrak{P}) := \dim A - \dim V(\mathfrak{P}) \), the codimension of \( V(\mathfrak{P}) \) in \( \text{Spec}(A) \) and if \( \mathfrak{p} \subset \mathfrak{P} \) is another prime ideal, hence \( V(\mathfrak{P}) \subset V(\mathfrak{p}) \), we define \( \text{codim}_{V(\mathfrak{p})} V(\mathfrak{P}) := \dim V(\mathfrak{p}) - \dim V(\mathfrak{P}) \), the codimension of \( V(\mathfrak{P}) \) in \( V(\mathfrak{p}) \). We wish to show in all cases that

(94.13) \( \text{ht} \mathfrak{p} = \text{codim}_{V(\mathfrak{p})} V(\mathfrak{P}) = \text{tr deg}_FA/\mathfrak{p} - \text{tr deg}_FA/\mathfrak{P} \).

In particular,

\[
\dim A/\mathfrak{P} + \text{ht} \mathfrak{P} = \dim A = \dim V(\mathfrak{P}) + \text{ht} \mathfrak{P}.
\]
Of course, if height is to be codimension, we would also want if \( p \subseteq \mathfrak{p} \) are prime ideals in \( A \) that

\[
\text{ht} \ \mathfrak{p} = \text{ht} \ p + \text{ht} \ \mathfrak{p} / p.
\]

Intuitively, still more should be expected. If \( \mathfrak{B} = (f_1, \ldots, f_n) \) is an ideal in \( A \), geometrically, we would want to view the \( f_i \) as hypersurfaces in \( V(\mathfrak{B}) \), e.g., if \( A = F[t_1, \ldots, t_N] \), then the \( f_i \) are polynomials in \( F[t_1, \ldots, t_N] \), and we can view this as the hypersurface \( f_i = 0 \) in \( F^N \).

Taking our cue from linear algebra, we would then expect if \( V(p) \) is an irreducible component of \( V(\mathfrak{B}) \) that

\[
\dim V(p) \geq \dim A - n,
\]

i.e., \( \text{codim}_{\text{Spec}(A)} V(p) \leq n \) or if height is the same as codimension that

\[
\text{ht} \ p \leq n.
\]

We shall, in fact, show all of this to be true for an affine \( F \)-algebra. Note that if \( F \) is algebraically closed, the above still will hold with \( Z(A) \) replacing \( V(\mathfrak{A}) \), etc. It will therefore follow that the algebraic notions of Krull dimension and height translate into the correct geometric notions.

Most of the above will follow by our next theorem. We first make the following definition:

**Definition 94.15.** A chain of prime ideals

\[
p_0 < p_1 < \cdots < p_m = p
\]

in a commutative ring is called **saturated** if no further new primes can be added to this chain.

We establish Equation 94.13.

**Theorem 94.16.** Let \( F \) be a field and \( A \) be an affine \( F \)-algebra that is a domain, \( p \) a prime ideal in \( A \), and

\[
0 \leq p_1 < \cdots < p_m = p
\]

a saturated chain of prime ideals in \( A \). Then

1. \( \text{ht} \ p = m \).
2. \( \dim A = \dim A / p + \text{ht} \ p \).

Moreover, \( \text{ht} \ p = \text{codim}_{\text{Spec}(A)} V(p) \) and all maximal saturated chains of prime ideals in \( p \) (hence starting from 0 and ending at \( p \)) have the same number of links. In particular, \( \dim A = \text{ht} m \) for all maximal ideals in \( A \).
Proof. Clearly, we need only prove (1) and (2). Let \( n = \dim A \), finite by Noether Normalization, and \( A = F[x_1, \ldots, x_r] \).

**Case 1.** If \( m = 1 \), then (1) and (2) hold.

As \( 0 < p \) is a saturate chain of prime ideals, \( \operatorname{ht} p = 1 \), so (1) holds.

**Subcase 1.** \( n = r \).

Since \( \dim A = \operatorname{tr} \deg_F A \), the elements \( x_1, \ldots, x_n \) are algebraically independent over \( F \). Therefore, \( A \cong F[t_1, \ldots, t_n] \) is a UFD. In particular, \( p \) must contain a prime element \( f \), so we must have \( \operatorname{ht} p = 1 \) and \( p = (f) \). Since \( f \) is nonzero, there exists an \( i \) such that \( x_i \) occurs non-trivially in \( f \), say \( i = n \). Let \( \pi : A \to A/p \) be the canonical epimorphism, so \( \pi(x_1, \ldots, x_n) = f(x_1, \ldots, x_n) = 0 \) in the domain \( \overline{A} = F[x_1, \ldots, x_n] \), i.e., \( x_1, \ldots, x_n \) are algebraically dependent over \( F \). Therefore, \( \operatorname{tr} \deg_F A = \dim F[x_1, \ldots, x_n] \), so \( \operatorname{tr} \deg_F A < n \). If \( g \in A \) satisfies \( \overline{g} = g(x_1, \ldots, x_n) = 0 \), then \( g \in \ker \pi = (f) \), so \( f \mid g \) in \( A \). In particular, \( g = 0 \) or \( x_n \) occurs non-trivially in \( g \). In particular, if \( h \in F[x_1, \ldots, x_{n-1}] \trianglelefteq F[t_1, \ldots, t_{n-1}] \) satisfies \( h(x_1, \ldots, x_{n-1}) = 0 \), then \( h = 0 \). It follows that \( x_1, \ldots, x_{n-1} \) are algebraically independent over \( F \), so \( \dim A = \dim F[x_1, \ldots, x_{n-1}] + 1 = \dim A + \operatorname{ht} p \).

**Subcase 2.** \( n \neq r \).

By Noether Normalization there exist \( y_1, \ldots, y_n \) in \( A \) algebraically independent over \( F \) satisfying \( A/F[y_1, \ldots, y_n] \) is integral. Let \( R = F[y_1, \ldots, y_n] \) and \( p' = p \cap R \), a prime ideal in \( R \). We have
\[
\dim A = \dim R \quad \text{and} \quad \operatorname{ht} p = \operatorname{ht} p',
\]
the second equality by Going Down as \( R \) is a normal domain. We also have the monomorphism \( R/p' \to A/p \) is integral, hence \( \dim R/p' = \dim A/p \). By Subcase 1 applied to \( R \), we have
\[
\dim A = \dim R = \dim R/p' + \operatorname{ht} p' = \dim A/p + \operatorname{ht} p = \dim(A/p) + 1,
\]
which is (2).

**Case 2.** Suppose that \( p \) is a maximal ideal in \( A \). Then \( \operatorname{ht} p = m = \dim A \).

We have a saturated chain
\[
0 < p_1 < \cdots < p_m = p.
\]
We show that \( m = \operatorname{ht} p = \dim A \) by induction on \( \dim A \). If \( \dim A = 1 \), the result is immediate, so we may assume that \( \dim A > 1 \). By Noether Normalization, there exist \( y_1, \ldots, y_n \) in \( A \) algebraically independent over \( F \) satisfying \( A/F[y_1, \ldots, y_n] \) is integral. Let \( R = F[y_1, \ldots, y_n] \), a normal domain, so by Going Down, we have \( \operatorname{ht} p = \operatorname{ht}(p \cap R) \). Moreover,
the inclusion \( R \hookrightarrow A \) is an integral map, so \( p \cap R \) is a maximal ideal in \( R \). Since \( \dim A = \dim R \) and
\[
0 < p_1 \cap R < \cdots < p_m \cap R = p \cap R
\]
by Comparability, it suffices to show that \( m = n \) in the case that \( A = R \). So assume that \( A = R \). We know, by the definition of \( \dim A \), that
\[
m \leq \dim A = \tr \deg_K A = n.
\]
By the argument in Case 1, Subcase 1, we have
\[
n - 1 = \dim A/p_1 = \tr \deg_K A - 1.
\]
Let \( - : A \to A/p_1 \), the canonical epimorphism. Then we have
\[
0 = \overline{p}_1 < \cdots < \overline{p}_m = \overline{p}
\]
with \( \overline{p} \) a maximal ideal in \( \overline{A} \) by the Correspondence Principle. By induction on \( n \), we have \( m - 1 = n - 1 \), so \( m = n \). This establishes Case 2.

**Case 3.** For any prime ideal \( p \) in \( A \), (1) and (2) hold.

Let \( m \in V(p) \cap \text{Max}(A) \). By Case 2, every saturated chain of prime ideals in \( m \) has \( n = \text{ht} m = \dim A \) links. Extend the saturated chain
\[
0 < p_1 < \cdots < p_m = p
\]
in \( A \) to a saturated chain of prime ideals
\[
0 < p_1 < \cdots < p_m < \cdots < p_n = m
\]
in \( A \) using the fact that \( A \) is Noetherian. Since \( m/p_m \) is a maximal ideal in \( A/p_m \), by Case 2, we have
\[
\text{ht} m/p_m = \dim A/p_m = n - m.
\]
Clearly,
\[
\dim A \geq \dim A/p_m + \text{ht} p_m \quad \text{and} \quad \text{ht} p_m \geq m,
\]
so
\[
\dim A \geq \dim A/p_m + \text{ht} p_m
\]
\[
= (n - m) + \text{ht} p_m \geq (n - m) + m = n = \dim A,
\]
hence
\[
\dim A = \dim A/p + \text{ht} p \quad \text{and} \quad \text{ht} p = m
\]
as needed. \( \square \)
A Noetherian ring $R$ is called *catenary* if all prime ideals $p < \mathfrak{p}$ in $R$ and saturated chains of prime ideals beginning at $p$ and ending at $\mathfrak{p}$ have the same finite number of links and is called *universally catenary* if every finitely generated commutative $R$-algebra is catenary. It is easy to show the following:

**Lemma 94.17.** If $R$ is a catenary ring and $\mathfrak{A} \subseteq R$ an ideal, then $R/\mathfrak{A}$ is a catenary ring.

As every affine $F$-algebra that is also a domain, e.g., $F[t_1, \ldots, t_n]$, is catenary by the theorem, by the lemma, we conclude that:

**Corollary 94.18.** Every field is universally catenary.

We turn to the last fact that we wish to prove, viz., if $\mathfrak{A} = (f_1, \ldots, f_n)$ is an ideal in an affine $F$-algebra, then $ht \mathfrak{p} \leq n$. We prove this in greater generality. This result, due to Krull, is one of the foundational theorems in commutative algebra.

Recall if $\mathfrak{A} \subseteq R$ is an ideal, the *height* of $\mathfrak{A}$ is

$$ht \mathfrak{A} := \inf \{ht \mathfrak{p} \mid \mathfrak{p} \in V(\mathfrak{A})\} = \inf \{ht \mathfrak{p} \mid \mathfrak{p} \in V(\mathfrak{A}) \text{ minimal}\}.$$

**Theorem 94.19.** (Principal Ideal Theorem) Let $R$ be a commutative Noetherian ring, $\mathfrak{A} \subseteq R$ an ideal and $\mathfrak{p} \in V(\mathfrak{A})$ minimal. If $\mathfrak{A} = (x_1, \ldots, x_n)$, then $ht \mathfrak{p} \leq n$. In particular, $ht \mathfrak{A} \leq n$.

**Proof.** The crucial result is the case that $n = 1$, hence the theorem's name.

**Case 1:** $\mathfrak{A} = (x)$ is principal.

If $x = 0$ or even nilpotent, then $\mathfrak{p} \in \text{Min}(R)$ and $ht \mathfrak{p} = 0$. So we may assume that $x$ is nonzero. Suppose that the result is false, i.e., $ht \mathfrak{p} > 1$. Then we have a chain of prime ideals in $R$,

$$\mathfrak{p}_0 < \mathfrak{p}_1 < \mathfrak{p} \text{ with } \mathfrak{p}_0, \mathfrak{p}_1 \not\in V(\mathfrak{A}),$$

since $\mathfrak{p} \in V(\mathfrak{A})$ is minimal. We make some reductions. By the Correspondence Principle, we have $\mathfrak{p}/\mathfrak{p}_0 \in V\left(\left((x) + \mathfrak{p}_0\right)/\mathfrak{p}_0\right)$ is minimal and $ht \mathfrak{p}/\mathfrak{p}_0 \geq 2$. Replacing $R$ by $R/\mathfrak{p}_0$, we may assume that $R$ is a domain and $\mathfrak{p}_0 = 0$. Next let $S = R \setminus \mathfrak{p}$. Then $R_\mathfrak{p} = S^{-1}R$ satisfies $S^{-1}\mathfrak{p} \in V(R_\mathfrak{p}x)$ is minimal and $0 < S^{-1}\mathfrak{p}_1 < S^{-1}\mathfrak{p}$. So replacing $R$ by $R_\mathfrak{p}$ and changing notation (considerably), we are reduced to $R$ satisfying the following conditions:

1. $R$ is a Noetherian domain
2. $(R, \mathfrak{m})$ is a local ring.
3. $ht \mathfrak{m} \geq 2$, so there exists a chain of prime ideals $0 < \mathfrak{p} < \mathfrak{m}$ in $R$. 


(4) The element $x \in m$ satisfies $V(x) = \{m\}$. In particular, $x / \mathfrak{P}$. Let $\mathfrak{A}_i = (\mathfrak{P}R_{\mathfrak{q}})^i \cap R < R$ a finitely generated ideal in $R$ for every non-negative integer $i$.

Note that $(R_{\mathfrak{q}}, \mathfrak{P}R_{\mathfrak{q}})$ is a Noetherian local domain with $R \subset R_{\mathfrak{q}} \subset qf(R)$.

We have a descending chain of ideals in $R$:

$$\mathfrak{P} = \mathfrak{A}_1 \supset \mathfrak{A}_2 \supset \cdots \supset \mathfrak{A}_i \supset \cdots.$$ 

As $m \in V(x)$ is minimal, $\text{Spec}(R/(x)) = \{m/(x)\}$, hence $R/(x)$ is a Noetherian ring of dimension zero. By Akizuki’s Theorem 93.9, $R/(x)$ is an Artinian ring. Therefore, the descending chain of ideals

$$(\mathfrak{A}_1 + (x))/(x) \supset \cdots \supset (\mathfrak{A}_i + (x))/(x) \supset \cdots$$

in $R/(x)$ stabilizes, say $(\mathfrak{A}_j + (x))/(x) = (\mathfrak{A}_{j+n} + (x))/(x)$ for all $n \geq 0$. By the Correspondence Theorem, $(\mathfrak{A}_j + (x)) = (\mathfrak{A}_{j+n} + (x))$.

Let $a \in \mathfrak{A}_j$. Then

$$a = b + rx, \text{ for some } b \in \mathfrak{A}_{j+1}, \text{ and } r \in R.$$ 

In particular, $rx = a - b \subset \mathfrak{A}_j \subset (\mathfrak{P}R_{\mathfrak{q}})^j \subset R_{\mathfrak{q}}$. As $x \notin \mathfrak{P}$, it must be a unit in $R_{\mathfrak{q}}$, consequently, $r \in (\mathfrak{P}R_{\mathfrak{q}})^j$. It follows that $r \in (\mathfrak{P}R_{\mathfrak{q}})^j \cap R = \mathfrak{A}_j$, i.e.,

$$\mathfrak{A}_j = \mathfrak{A}_{j+1} + \mathfrak{A}_j x \subset \mathfrak{A}_{j+1} + \mathfrak{A}_j m \subset \mathfrak{A}_j.$$ 

Therefore, $\mathfrak{A}_j = \mathfrak{A}_{j+1} + \mathfrak{A}_j m$ in the local ring $(R, m)$ with $\mathfrak{A}_j$ finitely generated. By alternate form of Nakayama’s Lemma, Corollary 90.11, this implies that we have $\mathfrak{A}_j = \mathfrak{A}_{j+1}$. Thus,

$$(\mathfrak{P}R_{\mathfrak{q}})^j = \mathfrak{A}_j R_{\mathfrak{q}} = \mathfrak{A}_{j+1} R_{\mathfrak{q}} = (\mathfrak{P}R_{\mathfrak{q}})^{j+1}$$

in the local Noetherian ring $(R_{\mathfrak{q}}, \mathfrak{P}R_{\mathfrak{q}})$. Hence $(\mathfrak{P}R_{\mathfrak{q}})^j = 0$ in the domain $R_{\mathfrak{q}}$, so $\mathfrak{P}R_{\mathfrak{q}} = 0$ and $\mathfrak{P} = 0$, a contradiction. This proves Case 1.

Case 2. General case (the proof is due to Akizuki).

Let $\mathfrak{A} = (a_1, \ldots, a_n) < R$ be an ideal with $p \in V(\mathfrak{A})$ minimal. Suppose that $\text{ht } p > n$. Choose a prime ideal $p_1$ in $R$ with $p_1 < p$, and $\text{ht } p_1 \geq n$. Replacing $R$ by $R_p$ if necessary, we may assume that $(R, m)$ is a local ring with $m = p$, so $V(\mathfrak{A}) = \{m\}$. As $R$ is Noetherian, we may assume that $p_1 < m$ is saturated, i.e., there exists no prime ideal $q$ in $R$ with $p_1 < q < m$. (Why?) Since $p_1 \notin V(\mathfrak{A})$, there exists an $i$ such that $a_i \notin p_1$. We may assume that $a_1 \notin p_1$. Therefore, we have

$$p_1 < p_1 + (a_1) \subset p = m.$$
Hence $V(p_1 + (a_1)) = \{m\}$, as $p_1 < m$ is saturated. Consequently, we have

\[
\sqrt{p_1 + (a_1)} = \bigcap_{\mathfrak{p} \in V(p_1 + (a_1))} \mathfrak{p} = m.
\]  

Since $m$ is finitely generated, there exists a positive integer $k$ satisfying $m^k \subset p_1 + (a_1)$ by (*) (and the definition of the radical). As we also have $\mathfrak{A} \subset m$, we conclude that $\mathfrak{A}^k \subset m^k \subset p_1 + (a_1)$. Therefore, we have equations

\[a_i^k = b_i + c_i a_1 \text{ for some } b_i \in p_i \text{ and } c_i \in R, \ i = 2, \ldots, n.\]

Let $\mathfrak{B} = (b_2, \ldots, b_n) \subset p_1$, so $p_1 \in V(\mathfrak{B})$. As $\text{ht} \ p_1 \geq n$, by induction, we conclude that there exists a prime ideal $q \in V(\mathfrak{B})$ minimal with $\mathfrak{B} \subset q < p_1$. However, $a_1^k, \ldots, a_n^k$ all lie in $\mathfrak{B} + (a_1) \subset q + (a_1)$, so

\[\{m\} = V(\mathfrak{A}) = V(a_1, \ldots, a_n) = V(a_1^k, \ldots, a_n^k) \supset V(q + (a_1)).\]

Therefore, $V(q + (a_1)) = \{m\}$ and $m/q \in V((q + (a_1))/q)$ is minimal (working in the ring $R/q$). By Case 1, we know that $\text{ht} \ m/q \leq 1$, which contradicts the existence of the chain of primes $0 = q/q < p_1/q < m/q$ in $R/q$. $\square$

The Principal Ideal Theorem allows us to show Equation 94.14 is valid.

**Corollary 94.20.** Let $F$ be a field and $A$ be an affine $F$-algebra that is a domain of dimension $n$. If $f_1, \ldots, f_m \in A$ with $m < n$, then $\dim V(f_1, \ldots, f_n) \geq n - m > 0$.

**Proof.** Let $p$ a minimal prime in $(f_1, \ldots, f_m)$. Then $\text{ht} \ p \leq m$ by the Principal Ideal Theorem. The result follows by Theorem 94.16. $\square$

Since ideals in a Noetherian ring are finitely generated, we have:

**Corollary 94.21.** Let $R$ be a Noetherian ring and $\mathfrak{p}$ a prime ideal in $R$. Then $\text{ht} \ \mathfrak{p}$ is finite. In particular, the set of prime ideals in $R$ satisfies the descending chain condition.

**Corollary 94.22.** Every local Noetherian ring has finite dimension.

**Corollary 94.23.** Let $R$ be a Noetherian domain. Then $R$ is a UFD if and only if every height one prime in $R$ is principal if and only if the set of height one primes is equal to the set of nonzero principal prime ideals.
We know by Kaplansky’s Theorem 30.1 that $R$ is a UFD if and only if every nonzero prime ideal in $R$ contains a prime element. As every nonzero prime ideal in $R$ contains a prime ideal of height one by Corollary 94.21, the result follows.

We shall also need a further generalization of the Principal Ideal Theorem.

**Theorem 94.24.** Let $R$ be a Noetherian ring, $\mathfrak{A} < R$ an ideal generated by $n$ elements, and $p \in V(\mathfrak{A})$ a prime satisfying $\text{ht}_{R/\mathfrak{A}} p/\mathfrak{A} = m$. Then

$$\text{ht}_R p \leq m + n = h_{R/\mathfrak{A}} p/\mathfrak{A} + n.$$  

**Proof.** We induct on $m$ (and all $n$ simultaneously).

$m = 0$: As $p \in V(\mathfrak{A})$ is minimal, this follows from the Principal Ideal Theorem.

$m > 1$: As $R$ is Noetherian, there exist finitely prime ideals minimal over $\mathfrak{A}$, say $p_1, \ldots, p_r$. As $p$ is not minimal, $p \not\subseteq \bigcap_{i=1}^r p_i$ by the Prime Avoidance Lemma 90.16. Let $a \in p \setminus \bigcap_{i=1}^r p_i$ and set $\mathfrak{B} = \mathfrak{A} + Ra$. Then $p_i \not\in V(\mathfrak{B})$ for $i = 1, \ldots, r$, so there exist no chain of prime ideals in $V(\mathfrak{B})$ starting at a $p_i$ and ending in $p$. It follows that $\text{ht}_{R/\mathfrak{B}} p \leq m - 1$, since $V(\mathfrak{B}) \subset V(\mathfrak{A})$. By induction, we have

$$\text{ht} p \leq \text{ht}_{R/\mathfrak{B}} p/\mathfrak{B} + (n + 1) \leq n + m.$$  

**Corollary 94.25.** Let $R$ be a Noetherian ring and $p$ a prime ideal in $R$ of height $m$. Suppose that $x$ is an element in $p$.

1. We have $m - 1 \leq \text{ht}_{R/(x)} p/(x) \leq m$.
2. Suppose that $x$ does not lie in any minimal prime of $R$. Then $\text{ht}_{R/(x)} p/(x) = m - 1$.
3. If $x$ is not a zero divisor in $R$, then $\text{ht}_{R/(x)} p/(x) = m - 1$.

**Proof.** Let $s = \text{ht}_{R/(x)} p/(x)$ with $p \in V(x)$. By the Principal Ideal Theorem 94.19, $m = \text{ht} p \leq s + 1$, and it is immediate that $s = \text{ht}_{R/(x)} p/(x) \leq \text{ht} p = m$. Therefore, (1) follows. Suppose that the set $\text{Min}(R) \cap V(x)$ is empty. Then $\text{ht}_{R/(x)} p/(x) < \text{ht} p$, so (2) follows. Finally, as $\bigcup_{\text{Min}(R)} p \subset \text{zd}(R)$ by Proposition 91.15, we have (2) implies (3).

**Remarks 94.26.** Let $R$ be a Noetherian ring.

1. Let $\mathfrak{A} = (a_1, \ldots, a_n) < R$ be of height $n$. Then there exist a minimal prime ideal $p \in V(\mathfrak{A})$ satisfying $\text{ht} p = \text{ht} \mathfrak{A} = n$ and $\mathfrak{A}$ cannot be generated by a fewer number of elements.
2. It is possible for an ideal \( \mathfrak{A} < R \) to have \( p_1, p_2 \in V(\mathfrak{A}) \) both minimal but of different heights. For example, let \( R = F[X, Y] \) with \( F \) a field and \( X, Y \) indeterminants. Set \( \mathfrak{A} = (X(X - 1), XY) \). Then

\[
\mathfrak{A} = (X) \cup (X - 1, Y), (X), (X - 1, Y) \in V(\mathfrak{A})
\]

where \( (X) \) and \( (X - 1, Y) \) are minimal primes of \( \mathfrak{A} \) of heights 1 and 2 respectively, and \( V(\mathfrak{A}) = V(X) \cup V(X - 1, Y) \), i.e., \( V(\mathfrak{A}) \) is the union of the \( Y \)-axis and the point \((1, 0)\) in \( F^2 \).

The next result can be viewed as a converse to the Principal Ideal Theorem.

**Corollary 94.27.** Let \( R \) be a Noetherian ring, \( \mathfrak{A} < R \) an ideal, and \( \mathfrak{p} < R \) a prime ideal.

1. If \( \text{ht} \mathfrak{A} = n \geq 1 \), then there exist \( a_1, \ldots, a_n \in \mathfrak{A} \) satisfying \( \text{ht}(a_1, \ldots, a_i) = i \) for \( i = 1, \ldots, n \).
2. If \( \text{ht} \mathfrak{p} = n \geq 1 \), then there exist \( a_1, \ldots, a_n \in \mathfrak{p} \) satisfying \( \mathfrak{p} \in V(a_1, \ldots, a_n) \) is minimal and \( \text{ht}(a_1, \ldots, a_i) = i \) for \( i = 1, \ldots, n \).

**Proof.** (1): As \( \text{ht} \mathfrak{A} \geq 1 \) and \( \text{Min}(R) \) is finite, there exists an element \( a_1 \in \mathfrak{A} \setminus \bigcup_{\mathfrak{p} \in \text{Min}(R)} \mathfrak{p} \) by the Prime Avoidance Lemma. Therefore, \( \mathfrak{p} \notin V(a_1) \) for all primes \( \mathfrak{p} \) in \( \text{Min}(R) \), so \( \text{ht}(a_1) > 0 \). By the Principal Ideal Theorem, \( \text{ht}(a_1) \leq 1 \), hence \( \text{ht}(a_1) = 1 \). By induction, there exists \( a_1, \ldots, a_i \in \mathfrak{A} \) satisfying \( \text{ht}(a_1, \ldots, a_i) = i < n \). Let \( \mathfrak{p} \in V(a_1, \ldots, a_i) \) be minimal. By the Principal Ideal Theorem, \( \text{ht} \mathfrak{p} \leq i \) and hence \( \text{ht} \mathfrak{p} = i \). As \( \{ \mathfrak{p} \mid \mathfrak{p} \in V(a_1, \ldots, a_i), \text{minimal} \} \) is a finite set, there exists an element

\[
a_{i+1} \in \mathfrak{A} \setminus \bigcup_{\mathfrak{p} \in V(a_1, \ldots, a_i), \text{minimal}} \mathfrak{p}
\]

by the Prime Avoidance Lemma. Let \( \mathfrak{B} = (a_1, \ldots, a_{i+1}) \) and \( \mathfrak{Q} \in V(\mathfrak{B}) \subset V(a_1, \ldots, a_i) \). Then there exists a prime ideal \( \mathfrak{P} \) in \( V(a_1, \ldots, a_i) \) minimal satisfying \( \mathfrak{P} \subset \mathfrak{Q} \). Consequently, \( \text{ht} \mathfrak{Q} > \text{ht} \mathfrak{P} = i \). By the Principal Ideal Theorem, \( \text{ht} \mathfrak{B} \leq i + 1 \), hence \( \text{ht} \mathfrak{B} = i + 1 \) and (1) follows by induction.

(2): Choose \( a_i \) as in (1) applied to \( \mathfrak{A} = \mathfrak{p} \). Since \( \mathfrak{p} \in V(A_1, \ldots, a_n) \) satisfies \( \text{ht} \mathfrak{p} = \text{ht}(a_1, \ldots, a_n) \), we have \( \mathfrak{p} \in V(a_1, \ldots, a_n) \) is minimal. \( \square \)

**Corollary 94.28.** Let \( R \) be an affine \( F \)-domain of dimension \( n \) and \( \mathfrak{P} \) a prime ideal in \( R \) of height \( m \geq 1 \). Then there exist \( a_1, \ldots, a_m \) in \( R \) such that \( \mathfrak{P} \) is minimal over \( (a_1, \ldots, a_m) \). In particular, if \( \mathfrak{p} \in \)
$V(a_1,\ldots,a_m)$ is minimal, then $\dim R/p \geq n - m$ with equality if $p = (a_1,\ldots,a_m)$.

**Proof.** The result follows first applying Theorem 94.16 then by localizing at $p$ and applying Corollary 94.27.

The last corollary implies if $F$ a field over an algebraically closed field (respectively any field) and $Z$ is an irreducible affine variety (respectively, abstract irreducible variety) of codimension $r$ in $F^n$, then $Z$ is an irreducible component of an affine variety (respectively abstract variety) defined as the intersection of $r$ hypersurfaces in $F^n$.

**Exercises 94.29.**

1. Prove the following generalization of the Normalization Theorem 94.1: Let $F$ be a field and $A$ an affine $F$-algebra (not necessarily a domain). If

$$\mathfrak{A}_1 < \cdots < \mathfrak{A}_m < A$$

is a chain of ideals in $A$, then there exists a non-negative integer $n$ and elements $y_1,\ldots,y_n$ in $A$ algebraically independent over $F$ and integers

$$0 \leq h_1 \leq \cdots \leq h_m \leq n$$

satisfying the following:

(a) $F[y_1,\ldots,y_n]$ is integral in $A$.

(b) $A$ is a finitely generated $F[y_1,\ldots,y_n]$-module.

(c) $\mathfrak{A}_i \cap F[y_1,\ldots,y_n] = (y_1,\ldots,y_{h_i})$ for $i = 1,\ldots,m$.

2. Prove Lemma 94.17.

3. Let $R$ be a catenary ring and $S$ a multiplicative set in $R$. Then $S^{-1}R$ is catenary.

[Note. If $A$ is an affine $F$-algebra, this says that localizations of it are catenary. Such a localization may not be an affine $F$-algebra.]

4. Let $R$ be a Noetherian ring of dimension greater than one. Prove that $\text{Spec}(R)$ is infinite.

**95. Regular Local Rings**

In geometry, the notion of a simple (or non-singular) point of an algebraic variety is a fundamental notion. It depends on the Jacobian matrix at the point being of the maximal rank (cf. Remark 95.31 below). Indeed one of the interesting aspects of algebraic geometry is that all points need not be simple. Unfortunately, the algebraic analogue does not have all the properties that such geometric points have. We shall study the analogue in this section. In geometry, if a point is simple then the ring of germs at the point (i.e., the functions
defined in a neighborhood of the point with two such being identified if they agree in some neighborhood of the point) is a UFD. [If all points are simple, this leads to the theory of Weil divisors.] This is also true of the algebraic analogue. Unfortunately, the general proof depends on developing sufficient homological algebra, which we do not do. Instead we shall only prove the algebraic case corresponding to the geometric case.

Recall if $R$ is a local ring with maximal ideal $m$, we write this as $(R, m)$ is a local ring.

**Definition 95.1.** Let $(R, m)$ be a Noetherian local ring of dimension $n$. By Corollary 94.27, there exist elements $x_1, \ldots, x_n$ in $m$ satisfying $\text{ht}(x_1, \ldots, x_n) = \text{ht} m$. Hence $m \in V(x_1, \ldots, x_n)$ is minimal. We say that the elements $x_1, \ldots, x_n$ form a system of parameters for $R$. As $V(x_1, \ldots, x_n) = \{m\}$, the ideal $m/(x_1, \ldots, x_n)$ in $R/(x_1, \ldots, x_n)$ is nilpotent. Therefore, we have $m^k \subset (x_1, \ldots, x_n) \subset m$ for some integer $k$, since $m$ is finitely generated. If there exists a system of parameters that generate $m$, such a system of parameters is called a regular system of parameters for $R$ and $R$ is called a regular local ring.

Define

$$V\text{-dim } R := \dim_{R/m} m/m^2.$$  

By Corollary 90.12 of Nakayama’s Lemma, a minimal generating set for $m$ is precisely one whose images give a basis for the $R/m$-vector space $m/m^2$. Therefore, $V\text{-dim } R$ is the size of a minimal generating set for $m$. By the Principal Ideal Theorem 94.19,

$$\dim R = \text{ht } m \leq V\text{-dim } R.$$  

Therefore, $R$ is a regular local ring if and only if a minimal generating set for $m$ consists of $\dim R$ elements if and only if there exists a regular system of parameters for $m$. Regular local rings will be the analogue of a simple point (more exactly, the germs of functions at a simple point).

**Example 95.2.** Let $F$ be a field and $R = F[[t_1, \ldots, t_n]]$, the formal power series in $t_1, \ldots, t_n$. Then $R$ is a local ring with maximal ideal $m = (t_1, \ldots, t_n)$ by Exercise 36.11(2) and is Noetherian by Example 29.13(2). Therefore, it is a regular local ring by Exercise 36.11(2). By Theorem 36.9, the regular local ring $R = F[[t_1, \ldots, t_n]]$ is a UFD.

We shall use this to prove that a regular local ring arising from a simple point of an affine variety in the geometric case is a UFD. More specifically, we shall show if $(R, m)$ is a regular local ring of dimension $n$ containing a copy of the field $F = R/m$, then $R$ embeds into $F[[t_1, \ldots, t_n]]$. Using the fact that $F[[t_1, \ldots, t_n]]$ is a UFD, we shall show that $R$ is also a UFD.


Our first goal is to show that a regular local ring is a domain. To do so we need another (important) theorem of Krull. [This is another form of the Krull Intersection Theorem, Exercise 39.12(5).] Recall that a nonzero polynomial \( f \in R[t_1, \ldots, t_n] \) over a commutative ring \( R \) is called homogeneous if every monomial in \( f \) has the same total degree. We say it is homogeneous of degree \( d \) if every monomial in \( F \) is of total degree \( d \), e.g., \( t_1^2 t_3^3 t_4 + t_3^2 t_4^3 \) is homogeneous of degree 6. [We also will say the zero polynomial is homogeneous.]

**Theorem 95.3. (Krull)** Let \( R \) be a Noetherian ring and \( \mathfrak{A} \) an ideal in \( R \). Suppose that \( 1+\mathfrak{A} \) consists of nonzero divisors. Then \( \bigcap_{k=0}^\infty \mathfrak{A}^k = 0 \).

**Proof.** Let \( \mathfrak{A} = (a_1, \ldots, a_n) \) and \( x \) an element of \( \bigcap_{k=0}^\infty \mathfrak{A}^k \). Let \( \mathfrak{a} = (a_1, \ldots, a_n) \in R^n \) and \( e_{\mathfrak{a}} : R[t_1, \ldots, t_n] \to R \) the evaluation at \( \mathfrak{a} \). Clearly, for each \( k > 0 \), there exists a homogeneous element \( f_k \) in \( R[t_1, \ldots, t_n] \) of degree \( k \) satisfying \( x = f_k(\mathfrak{a}) \). Let \( \mathfrak{B} \) be the ideal in \( R[t_1, \ldots, t_n] \) generated by the \( f_k, k > 0 \). Since \( R[t_1, \ldots, t_n] \) is a Noetherian ring by the Hilbert Basis Theorem 40.1, it satisfies ACC, so there exist \( f_1, \ldots, f_m \) and \( x = f_i(\mathfrak{a}) \) for \( i = 1, \ldots, m \) generating \( \mathfrak{B} \). In particular,

\[
f_{m+1} = g_1 f_1 + \cdots + g_m f_m,
\]

with \( g_i \in R[t_1, \ldots, t_n] \) homogeneous of degree \( m + 1 - i \geq 1 \). It follows that \( g_i(\mathfrak{a}) \) lies in \( \mathfrak{A}^{m+1-i} \subset \mathfrak{A} \) for \( i = 1, \ldots, m \), and

\[
x = f_{m+1}(\mathfrak{a}) = \sum_{i=1}^m g_i(\mathfrak{a}) f_i(\mathfrak{a}) = (\sum_{i=1}^m g_i(\mathfrak{a})) x.
\]

Set \( y = \sum_{i=1}^m g_i(\mathfrak{a}) \). Then \( (1-y)x = 0 \), and by hypothesis, \( x = 0 \). \( \square \)

**Corollary 95.4.** Let \( R \) be a Noetherian domain and \( \mathfrak{A} < R \) an ideal. Then \( \bigcap_{k=0}^\infty \mathfrak{A}^k = 0 \).

**Corollary 95.5.** Let \( R \) be a Noetherian ring and \( \mathfrak{A} < R \) an ideal lying in the Jacobson radical \( \text{rad}(R) \) of \( R \). Then for every ideal \( \mathfrak{B} \subset R \), we have \( \bigcap_{k=0}^\infty (\mathfrak{B} + \mathfrak{A}^k) = \mathfrak{B} \).

**Proof.** Let \( \sim : R \to R/\mathfrak{B} \) be the canonical ring epimorphism. Then \( \overline{\mathfrak{A}}^k = (\mathfrak{A}^k + \mathfrak{B})/\mathfrak{B} \). For all \( a \in \mathfrak{A} \subset \text{rad}(R) \), so \( 1 + a \in R^\times \) by Exercise 27.19(8). Therefore, \( 1 + \overline{a} \) in a nonzero divisor in \( R \) for all \( a \in \mathfrak{A} \). Hence

\[
\bigcap_{k=0}^\infty \overline{\mathfrak{A}}^k = \bigcap_{k=0}^\infty \left( \frac{\mathfrak{B} + \mathfrak{A}^k}{\mathfrak{B}} \right) = \overline{0}.
\]

The result follows. \( \square \)
Corollary 95.6. Let $R$ be a local Noetherian ring with maximal ideal $m$. Then $\bigcap_{k=0}^{\infty} m^k = 0$.

We need the corollary to prove the following lemma:

Lemma 95.7. Let $R$ be a local Noetherian ring. Suppose that $R$ is not a domain and $p = (p)$ is a prime ideal. Then $p$ is a minimal prime of $R$.

Proof. By Corollary 95.6, we have $\bigcap_{k=0}^{\infty} p^k = 0$. As $R$ is not a domain, it suffices to show if $q$ is a prime ideal in $R$ satisfying $q < p$, then $q \subset \bigcap_{k=0}^{\infty} p^k$. But if $q \in q$, then $q = pq_1$ in $p$. As $p \notin q$, $q_1 \in q$. Iterating this shows that $q \in \bigcap_{k=0}^{\infty} p^k$. □

Lemma 95.8. Let $R$ be a regular local ring of dimension $n$ with maximal ideal $m$. Let $x \in m \setminus m^2$. Then $R/(x)$ is a regular local ring of dimension $n - 1$.

Proof. Let $\varphi : R \to R/(x)$ be the canonical epimorphism. Let $x_1, \ldots, x_m$ lie in $m$ be chosen such that $x_1, \ldots, x_m$ is a minimal generating set for $m$.

We show that $x, x_1, \ldots, x_m$ is a minimal generating set for $m$. Certainly they generate, so suppose that we have a $R$-linear combination $ax + r_1 x_1 + \cdots + r_m x_m$ that lies in $m^2$. We must show each coefficient lies in $m$. As $\overline{r_1 x_1} + \cdots + \overline{r_m x_m}$ lies in $\overline{m}^2$, each $\overline{r_i}$ lies in $\overline{m}$ by the minimality of $x_1, \ldots, x_m$ (cf. Corollary 90.12 of Nakayama’s Lemma). It follows that $r_1, \ldots, r_m$ all lie in $m$, hence $ax \in m^2$. Since $x \notin m^2$, $a \in m$. Thus $x, x_1, \ldots, x_m$ is a minimal generating set. This means $V\text{-dim } R = m + 1$.

So to finish, we must show that $\text{ht } \overline{m} = n - 1$. We know that

$$\text{dim } R/(x) \leq (V\text{-dim } R) - 1 = (\text{dim } R) - 1 \leq \text{dim } R/(x),$$

and the result follows by Corollary 94.25. □

Corollary 95.9. Let $(R, m)$ be a Noetherian local ring with $x \in m \setminus m^2$ and lying in no minimal prime ideal of $R$. Then $R$ is a regular local ring if and only if $R/(x)$ is a regular local ring.

Proof. By the Lemma, we need only show if $R/(x)$ is regular, then $R$ is. But this follows from Corollary 94.25. □

We can now attain our goal.

Theorem 95.10. Let $R$ be a regular local ring. Then $R$ is a domain.

Proof. Let $R$ be of dimension $n$. If $n = 0$, then $R$ is a field so a domain, so we may assume that $n > 0$. Let $m$ be the maximal ideal
of \( R \) and nonzero. By Nakayama’s Lemma, \( m \neq m^2 \). Let \( x \in m \setminus m^2 \). Then \( R/(x) \) is a regular local ring of dimension \( n - 1 \) by Lemma 95.8, so a domain by induction. Suppose that \( R \) is not a domain. Then by Lemma 95.7, \( (x) \) is a minimal prime ideal. Since this is true for all \( x \in m \setminus m^2 \), we have \( m \setminus m^2 \subseteq \bigcup_{p \in \text{Min}(R)} p \). Since \( \text{Min}(R) \) is a finite set as \( R \) is Noetherian, \( m \subseteq p \), for some \( p \in \text{Min}(R) \) by the Prime Avoidance Lemma 90.16. But this means \( R \) is of dimension 0, a contradiction. \( \square \)

**Corollary 95.11.** Let \( R \) be a regular local ring and \( x \in m \setminus m^2 \). Then \( (x) \) is a prime ideal in \( R \). In particular, if \( x_1, \ldots, x_n \) is a regular system of parameters for \( R \), then \( R/(x_1, \ldots, x_i) \) is a regular local ring of dimension \( n - i \) and \( (x_1, \ldots, x_i) \) is a prime ideal of height \( i \) for \( i = 0, \ldots, n \).

**Remark 95.12.** It is a fact that if \( R \) is a regular local ring and \( p \) a prime ideal in \( R \) that \( R_p \) is a regular local ring. The proof uses homological algebra methods used to prove that \( R \) is shown to be a UFD for an arbitrary regular local ring.

**Remark 95.13.** We have seen that \( R \) is a regular local ring if and only if its maximal ideal is generated by a system of parameters. We comment on the generalization of this property. We omit proofs. Let \( R \) be a Noetherian ring. In 92.5, we have seen that if \( M \) is an \( R \)-module, then \( \text{zd}(M) = \bigcup_{p \in \text{Ass}_R(M)} p \) and understand the minimal primes in \( \text{Ass}_R(M) \). The embedded primes do have an effect. Suppose that \( R \) is also a local ring. We constructed a system of parameters \( x_1, \ldots, x_n \) for \( R \) based on the condition that \( x_i \) is not in any minimal prime containing \( (x_1, \ldots, x_{i-1}) \) for \( i = 1, \ldots, n \). This does not involve all of the associated primes of \( (x_1, \ldots, x_n) \), i.e., all the zero divisors of \( R/(x_1, \ldots, x_n) \). It turns out that all the zero divisors, i.e., all elements of the associated primes of \( R/(x_1, \ldots, x_n) \) are important.

A more general approach would be the following: Let \( R \) be a commutative ring and \( M \) an \( R \)-module. An \( R \)-sequence on \( M \) is an ordered sequence \( x_1, \ldots, x_n \) of elements in \( R \) satisfying:

1. The ideal \( (x_1, \ldots, x_n)M \subset M \).
2. The image of \( x_i \) in \( \text{zd} \left( M/(x_1, \ldots, x_{i-1})M \right) \) under the natural epimorphism is not a zero divisor for \( i = 1, \ldots, n \).

If \( M = R \), we call an \( R \)-sequence \( x_1, \ldots, x_n \) on \( R \) just an \( R \)-sequence. The first condition is included so that all the quotients \( M/(x_1, \ldots, x_i)M \) are not zero. As an example, if \( R_0 \) is a commutative ring, then \( t_1, \ldots, t_n \) is an \( R \)-sequence in \( R = R_0[t_1, \ldots, t_n] \).

In general, if \( M \) is an \( R \)-module and \( x_1, \ldots, x_n \) is an \( R \)-sequence on \( M \), then a permutation of this sequence may not be an \( R \)-sequence.
on $M$. If $R$ is a Noetherian ring and $M$ a nonzero finitely generated $R$-module, then it can be shown that any permutation of $x_1, \ldots, x_n$ is still an $R$-sequence on $M$ if $x_1, \ldots, x_n$ lie in the Jacobson radical of $R$. [In general, the problem will arise if the image of $x_i$ is a zero divisor in $zd\left(M/(x_1, \ldots, x_{i-1})M\right)$ for some $i$.] In particular, if $R$ is a local Noetherian ring and $M$ a nonzero finitely generated $R$-module, then any permutation of an $R$-sequence on $M$ is an $R$-sequence on $M$; and it can also be shown that, in this case, all maximal $R$-sequences on $M$ have the same length. In particular, in this case the maximal length of an $R$-sequence $M$ is well-defined and called the depth of $M$. [In general, if $R$ is any Noetherian ring with $A$ an ideal in $R$ and $M$ a finitely generated module with $M \neq A M$, then one can show that any two maximal $R$-sequences on $M$ lying in $A$ have the same finite length.] For example, suppose that $R$ is a regular local ring and $x_1, \ldots, x_n$ is a regular system of parameters. Then $R/(x_1, \ldots, x_i)$ is a domain for $i = 1, \ldots, n$. It follows that $x_1, \ldots, x_n$ is an $R$-sequence. In particular, depth $R = \dim R$. In general, depth $R \leq \dim R$. A Noetherian local ring is called a Cohen-Macaulay ring if depth $R = \dim R$. [A Noetherian ring is called a Cohen-Macaulay ring if $R_m$ is Cohen-Macaulay ring for every maximal ideal $m \in \text{Spec}(R)$.] Being a Cohen-Macaulay ring is weaker than being a regular local ring. In fact, for a local domain to be a Cohen-Macaulay ring is weaker than being a normal domain. For example, the ring $\mathbb{C}[X,Y]/(Y-X^2(X+1))$ giving the nodal curve of Remark 94.2 localized at $(X,Y)/(Y-X^2(X+1))$ (the ring at the origin of the curve in $\mathbb{C}^2$) can be shown to be Cohen-Macaulay but not normal, as the images of $X, Y$ is checked to give a regular sequence. One can also show that a Noetherian local ring is a Cohen-Macaulay if and only if every system of parameters is an $R$-sequence. The notion of being Cohen-Macaulay is quite useful in algebraic geometry as well as in commutative algebra.

We next turn to another useful characterization of a regular local ring. We shall use dehomogenization of a homogeneous polynomial, e.g., substitute $t/t, t_1/t, \ldots, t_n/t$ for the variables in such an $f$ thereby obtaining a polynomial in $n$ variables. Note if $R$ is a domain, $f \in R[t, t_1, \ldots, t_n]$ a homogeneous of degree $d$, and $f(x, x_1, \ldots, x_n) = 0$ with $x, x_1, \ldots, x_n \in R$ nonzero if and only if $x^N f(1, x_1/x, \ldots, x_n/x) = 0$ in $R$ for all $N \geq d$. We use this to prove the following lemma:

**Lemma 95.14.** Let $(R, m)$ be a Noetherian local domain of dimension $n + 1$. Suppose that $m \in V(x, x_1, \ldots, x_n)$ is minimal (e.g., $R$ is a
regular local ring) and \( f \in R[t, t_1, \ldots, t_n] \) is a homogeneous polynomial. If \( f(x, x_1, \ldots, x_n) = 0 \), then all coefficients of \( f \) lie in \( m \).

**Proof.** Set \( \mathfrak{P} = m R[t_1, \ldots, t_n] \), the extension of \( m \) to \( R[t_1, \ldots, t_n] \). As \( R[t_1, \ldots, t_n]/\mathfrak{P} \cong (R/m)[t_1, \ldots, t_n] \) is a domain, \( \mathfrak{P} \) is a prime ideal in \( R[t_1, \ldots, t_n] \) and minimal in \( V_{R[t_1, \ldots, t_n]}(x, x_1, \ldots, x_n) \). In particular, \( \dim \mathfrak{P} \leq n + 1 \) by the Principal Ideal Theorem 94.19. As any proper chain of primes for \( m \) gives a proper chain of primes for \( \mathfrak{P} \), we have \( \dim \mathfrak{P} = \dim R = n + 1 \).

Let \( \varphi : R[t_1, \ldots, t_n] \to R[x_1/x, \ldots, x_n/x] \) be the \( R \)-algebra epimorphism determined by \( t_i \mapsto x_i/x \), which lies in the quotient field of \( R \). To prove the result, it suffices to show \( \ker \varphi \subset \mathfrak{P} \).

We show that \( x \notin \mathfrak{p} \). If \( x \in \mathfrak{p} \), then \( x, x_1, \ldots, x_n \in \mathfrak{p} \), and we have \( (x, x_1, \ldots, x_n) \subset \mathfrak{p} \cap R \). Thus \( \mathfrak{p} \cap R = m \) and \( \mathfrak{P} \subset \mathfrak{p} \), a contradiction. Therefore, \( x \notin \mathfrak{p} \).

Now suppose that \( g \in \ker \varphi \). In particular, if \( k >> 0 \), then \( x^k g(x_1, \ldots, x_n) = 0 \). Let \( R_x \) denote the domain obtained by localizing \( R \) at \( S = \{ x^n \mid n \geq 0 \} \). Applying the Remainder Theorem (of the Division Algorithm) to the monic polynomials \( t_i - x_i/x \) to \( g \) sequentially, we get polynomials \( h_i \) in \( R_x[t_i, \ldots, t_n] \), for \( i = 1, \ldots, n \), and equations in \( R_x[t_1, \ldots, t_n] \):

\[
g = (t_1 - x_1/x)h_1 + g(x_1/x, t_2, \ldots, t_n) \\
\vdots \\
g(x_1/x, \ldots, x_{i-1}/x, t_i, \ldots, t_n) = (t_i - x_i/x)h_i + g(x_1/x, \ldots, x_{i-1}/x, t_{i+1}, \ldots, t_n) \\
\vdots \\
g(x_1/x, \ldots, x_{n-1}/x, t_n) = (t_n - x_n/x)h_n + g(x_1/x, \ldots, x_{n-1}/x).
\]

Lifting these equations back to \( R[t_1, \ldots, t_n] \), i.e., clearing all denominators of all these polynomials, we see that there exits \( k > 0 \) such that \( x^k g \in \mathfrak{P} \). Therefore, \( \ker \varphi \subset \mathfrak{P} \) as needed. \( \square \)

**Definition 95.15.** A commutative ring \( R \) will be called a graded ring (on index set \( \mathbb{N} = \{ i \in \mathbb{Z} \mid i \geq 0 \} \)) if \( R = \bigoplus_{i=0}^{\infty} R_i \) as additive groups with multiplication satisfying \( R_i R_j \subset R_{i+j} \) for all \( i, j \geq 0 \). (Note in this case, \( R_0 \) is a commutative ring.) An element \( r \) in \( R_q \) is called a homogeneous element of degree \( d \). A homomorphism of graded rings is a ring homomorphism \( \varphi : \bigoplus_{i \geq 0} R_i \to \bigoplus_{i \geq 0} S_i \) that preserves degree,
i.e., \( \varphi(R_i) \subset S_i \) for all \( i \geq 0 \) and, a graded ring isomorphism if, in addition, \( \varphi \) is also a ring isomorphism.

Examples 95.16. Let \( R_0 \) be a commutative ring.

1. The polynomial ring \( R = R_0[t_1, \ldots, t_n] \) is a graded ring with \( R_d \) the group generated by the monomials of total degree \( d \).
2. \( R = R_0[[t_1, \ldots, t_n]] \) is a graded ring with \( R_d \) the group generated by the monomials of total degree \( d \). Note that the natural map \( R_0[t_1, \ldots, t_n] \to R_0[[t_1, \ldots, t_n]] \) induced by \( t_i \mapsto t_i \) is a injective graded ring homomorphism.
3. Let \( R \) be a commutative ring and \( \mathfrak{A} \) an ideal in \( R \). The graded ring of \( R \) determined by \( \mathfrak{A} \), called the associated graded ring relative to \( \mathfrak{A} \), is defined to be

\[
G_{\mathfrak{A}}(R) = \bigoplus_{i=0}^{\infty} \mathfrak{A}^i / \mathfrak{A}^{i+1}
\]

where \( \mathfrak{A}^0 = R \) and the multiplication is induced by

\[
(x + \mathfrak{A}^{i+1})(y + \mathfrak{A}^{j+1}) := xy + \mathfrak{A}^{i+j+1} \text{ for all } x \in \mathfrak{A}^i, y \in \mathfrak{A}^j.
\]

The elements in \( \mathfrak{A}^d / \mathfrak{A}^{d+1} \) are the homogeneous elements of \( G_{\mathfrak{A}}(R) \) of degree \( d \).

Using Lemma 95.14, we have the following characterization of regular local rings.

Theorem 95.17. Let \( (R, \mathfrak{m}) \) be a Noetherian local ring of dimension \( n \). Then the following are equivalent.

1. The ring \( (R, \mathfrak{m}) \) is a regular local ring.
2. The rings \( (R/\mathfrak{m})[t_1, \ldots, t_n] \) and \( G_{\mathfrak{m}}(R) \) are isomorphic graded rings.
3. The vector space \( \mathfrak{m}/\mathfrak{m}^2 \) over the field \( R/\mathfrak{m} \) is \( n \)-dimensional.

Proof. (1) \( \Rightarrow \) (2): As \( R \) is a regular local ring, \( \mathfrak{m} = (x_1, \ldots, x_n) \) for some \( x_1, \ldots, x_n \). The canonical ring epimorphism \( - : R \to R/\mathfrak{m} \) induces the natural graded ring epimorphism \( R[t_1, \ldots, t_n] \to (R/\mathfrak{m})[t_1, \ldots, t_n] \) and the graded ring epimorphism \( \varphi : (R/\mathfrak{m})[t_1, \ldots, t_n] \to G_{\mathfrak{m}}(R) \) induced by \( t_i \mapsto x_i + \mathfrak{m}^2 \). Then \( \varphi \) is injective by Lemma 95.14.

(2) \( \Rightarrow \) (3): This is clear.

(3) \( \Rightarrow \) (1): By (3), a minimal generating set for \( \mathfrak{m} \) consists of \( n = \dim R \) elements, so is a regular system of parameters. Therefore, \( R \) is regular. \( \square \)
We next turn to our main goal, to show that a regular local ring in the geometric case is a UFD. The algebraic condition we need is the following:

**Definition 95.18.** Let \((R, \mathfrak{m})\) be a Noetherian local ring. We say that \(R\) is *equicharacteristic* if \(R\) contains an isomorphic copy of the field \(R/\mathfrak{m}\). For convenience of notation, we shall assume that \(F = R/\mathfrak{m}\) lies in \(R\) if this is the case.

We wish to show that under the hypothesis of an equicharacteristic Noetherian local ring, we can assign a power series to each element of \(R\). In general, such an assignment will not be unique. However, for a regular local ring, it shall be as we shall see.

The idea is the following: We have defined a graded ring isomorphism \(G_{\mathfrak{m}}(R) \rightarrow (R/\mathfrak{m})[t_1, \ldots, t_n]\) extending the canonical map \(R \rightarrow R/\mathfrak{m}\) and sending \(x_i + \mathfrak{m}^2 \mapsto t_i\). In particular, this defines an \(R/\mathfrak{m}\)-isomorphism \(\mathfrak{m}^d/\mathfrak{m}^{d+1} \rightarrow (R/\mathfrak{m})[t_1, \ldots, t_n]_d\). Since \(F \subset R\), this map fixes \(F\). We also have the natural graded ring monomorphism \((R/\mathfrak{m})[t_1, \ldots, t_n] \rightarrow F[[t_1, \ldots, t_n]]\). As power series can be truncated to give polynomials of various degrees, we wish to mimic the construction of Taylor series as in calculus. We then will show that the construction is unique when \(R\) is a regular local ring. For example, if \(R\) is of dimension one and \(x \in \mathfrak{m}\) then we would expect that the inverse of the unit \(1 - x\) should be assigned the power series \(\sum_{i=0}^{\infty} t^i\) corresponding to \((1, 1 + \mathfrak{m}, 1 + \mathfrak{m}^2, \ldots)\) with truncations lying in \(G_{\mathfrak{m}}(R)\). We first do the construction.

**Construction 95.19.** Let \((R, \mathfrak{m})\) be an equicharacteristic local Noetherian ring with \(F = R/\mathfrak{m}\) and \(- : R \rightarrow F\) the canonical epimorphism. We assume that \(F\) also lies in \(R\) and \(-\) fixes \(F\). Suppose that \(\mathfrak{m} = (u_1, \ldots, u_m)\).

Let \(x \in R\). Define a power series in \(F[[t_1, \ldots, t_m]]\) associated to \(x\) as follows. Write

\[
x = \alpha_0 + x_1, \quad \text{with } \alpha_0 \in F, x_1 \in \mathfrak{m} \quad \text{and set } \overline{x} = \alpha_0.
\]

Then we have

\[
x_1 = \sum_{i=1}^{m} \alpha_i u_i + x_2 \quad \text{with } \alpha_i \in F, x_2 \in \mathfrak{m}^2.
\]

Similarly, define

\[
x_2 = \sum_{i,j=1}^{m} \alpha_{ij} u_i u_j + x_3 \quad \text{with } \alpha_{ij} \in F, x_3 \in \mathfrak{m}^3.
\]
Continuing in this way we see that for each \( r \), we can find homogeneous polynomials \( f_i \in F[t_1, \ldots, t_m] \) of degree \( i \), \( i = 0, \ldots, r \), satisfying
\[
x = \sum_{i=0}^{r} f_i(u_1, \ldots, u_m) + x_{r+1} \quad \text{with} \quad x_{r+1} \in \mathfrak{m}^{r+1}.
\]
Thus we can assign to \( x \) a formal power series
\[
\Phi(x) = f_0 + f_2 + f_3 + \cdots \quad \text{in} \quad F[[t_1, \ldots, t_m]].
\]

**Definition 95.20.** Let \((R, \mathfrak{m})\) be an equicharacteristic local Noetherian ring with maximal ideal \( \mathfrak{m} = (u_1, \ldots, u_m) \). If \( x \in R \), then a formal power series \( \Phi(x) = f_0 + f_2 + f_3 + \cdots \in F[[t_1, \ldots, t_m]] \) is called a (formal) Taylor series of \( x \) relative to \( u_1, \ldots, u_m \) if for all \( r \geq 0 \),
\[
x - \sum_{i=0}^{r} f_i(u_1, \ldots, u_m) \quad \text{lies in} \quad \mathfrak{m}^{r+1}.
\]

**Example 95.21.** Let \( F \) be a field and \( R = F[t]((t)) \), so \( \mathfrak{m} = (t) \). Let \( f = p/q \) with \( p, q \in F[t] \) and \( q(0) \neq 0 \). Then \( f \in R \). The Taylor series relative to \( t \) is a formal power series \( \sum_{i=0}^{\infty} \alpha_i t^i \) such that \( (p/q) - \sum_{i=0}^{r} \alpha_i t^i \equiv 0 \mod t^{r+1} \) for all \( r \geq 0 \). For example, \( 1/(1-t) \) has a Taylor series \( \sum_{i=0}^{\infty} t^i \), since
\[
\frac{1}{1-t} - \sum_{i=0}^{r} t^i = \frac{t^{r+1}}{1-t} \equiv 0 \mod t^{r+1} \quad \text{for all} \quad r \geq 0.
\]

**Proposition 95.22.** Let \((R, \mathfrak{m})\) be an equicharacteristic regular local ring with \( x_1, \ldots, x_n \) a regular system of parameters. Then for all \( x \in R \), there exist a unique Taylor series for \( x \) relative to \( x_1, \ldots, x_n \).

**Proof.** If \( \Phi(x) = f_0 + f_1 + \cdots \) and \( \Phi'(x) = f_0' + f_1' + \cdots \) are two Taylor series for \( x, x \in R \), then 0 is a Taylor series for \( \Phi(x) - \Phi'(x) \). So it suffices to show if 0 has a Taylor series \( \Phi = f_0 + f_1 + \cdots \), then \( \Phi = 0 \).

Suppose that \( f_r \) is the first nonzero homogenous term in \( \Phi \). Then \( f_r(x_1, \ldots, x_n) \) lies in \( \mathfrak{m}^{r+1} \). As a regular local ring is a domain, all coefficients of \( f_r \) lie in \( \mathfrak{m} \) by Lemma 95.14, so \( f_r(x_1, \ldots, x_n) = 0 \). \( \square \)

The proposition implies:

**Corollary 95.23.** Let \((R, \mathfrak{m})\) be an equicharacteristic regular local ring of dimension \( n \) with \( F = R/\mathfrak{m} \). If \( \mathfrak{m} = (x_1, \ldots, x_n) \), then there exists an \( F \)-algebra monomorphism \( R \rightarrow F[[t_1, \ldots, t_n]] \).
We now turn to proving that an equicharacteristic regular local ring is a UFD. We need some preliminaries. It is useful to set up the following notation.

**Notation 95.24.** Let \((R, \mathfrak{m})\) be an equicharacteristic regular local ring of dimension \(n\) with \(F = R/\mathfrak{m}\). As before we view \(F \subset R\) and as constructed before, we have an \(F\)-algebra monomorphism \(R \to F[[t, \ldots, t_n]]\) by \(x \mapsto \Phi(x)\). We shall view this map as an inclusion. In particular, we shall identify \(x \in R\) with its Taylor series, \(\Phi(x) = \sum_{i=0}^{\infty} f_i \in F[[t, \ldots, t_n]]\), where \(f_i \in F[t_1, \ldots, t_n]\) is homogeneous of degree \(i\) relative to \(t_1, \ldots, t_n\).

We set \(\widehat{R} = F[[t_1, \ldots, t_n]]\) and if \(\mathfrak{m} = (x_1, \ldots, x_n)\), we identify \(x_i\) in \(R\) with \(t_i\) in \(\widehat{R}\). We shall also write \(x = (x_1, \ldots, x_n)\) when viewed as a point in \(F^n\).

We know that \(\widehat{R}\) is a regular local ring with maximal ideal \(\widehat{\mathfrak{m}} = (t_1, \ldots, t_n)\) and \(R\) and \(\widehat{R}\) are domains.

We begin by looking at a relationship between ideals in \(R\) and \(\widehat{R}\).

**Lemma 95.25.** In the Notation 95.24, for every \(i \geq 1\), we have \(\widehat{\mathfrak{m}}^i \cap R = \mathfrak{m}^i\).

**Proof.** Certainly, for every \(i \geq 1\), we have \(\mathfrak{m}^i \subset \widehat{\mathfrak{m}}^i \cap R\), so we need only show that \(\widehat{\mathfrak{m}}^i \cap R \subset \mathfrak{m}^i\). Let \(x \in R\) have Taylor series relative to \(x_1, \ldots, x_n\) given by \(\Phi(x) = \sum_{i=0}^{\infty} f_i\) with each \(f_j = 0\) for \(j < i\). Then by definition,

\[
x - \sum_{j=0}^{i-1} f_j(x) \equiv 0 \pmod{\mathfrak{m}^i}.
\]

as needed. \(\square\)

**Lemma 95.26.** In the Notation 95.24, if \(\mathfrak{A} \subset R\) is an ideal, then we have \(\widehat{R} \mathfrak{A} \cap R = \mathfrak{A}\).

**Proof.** Again it suffices to show \(\widehat{R} \mathfrak{A} \cap R \subset \mathfrak{A}\). As \(\mathfrak{A}\) is Noetherian, \(\mathfrak{A} = (a_1, \ldots, a_m)\) for some \(a_1, \ldots, a_m \in \mathfrak{A}\). Let \(x \in \widehat{R} \mathfrak{A} \cap R\). We can write

\[
x = g^{(1)}(x)a_1 + \cdots + g^{(m)}(x)a_m \quad \text{with} \quad g^{(i)} \in \widehat{R}
\]

and

\[
g^{(i)} = g_0^{(i)} + \cdots + g_n^{(i)} + \cdots \quad \text{the corresponding power series}.
\]
Each \( g_n^{(i)}(x) \) can be approximated up to an element in \( \hat{m}^{n+1} \) for each \( n \), so we can find \( h_n^{(i)}(x) \) in \( F[t_1, \ldots, t_m] \) with \( h_n^{(i)}(x) \in R \) and
\[
h_n^{(i)}(x) = g_0^{(i)}(x) + \cdots + g_n^{(i)}(x) \pmod{\hat{m}^{n+1}}.
\]
Therefore, we have
\[
\sum_{i=1}^{m} g_n^{(i)}(x) a_i = \sum_{i=1}^{m} \left( h_n^{(i)}(x) + g_i^{(i)}(x) - h_n^{(i)}(x) \right) a_i
\]
with \( \sum_{i=1}^{m} \left( g_n^{(i)}(x) - h_n^{(i)}(x) \right) a_i \) lying in \( \hat{m}^{n+1} \cap R = m^{n+1} \) by the last lemma. Consequently, \( \sum_{i=1}^{m} g_n^{(i)}(x) a_i \) lies in \( \mathfrak{A} + m^{n+1} \) for every \( n \geq 1 \). Since \( \bigcap_{i=0}^{\infty} (\mathfrak{A} + m^n) = \mathfrak{A} \) by Corollary 95.5, the lemma is proven. \( \square \)

To prove that an equicharacteristic regular local ring \( R \) is a UFD, we only need to show that every irreducible element in \( R \) is a prime element, as \( R \) is Noetherian. We want to use the fact that \( \hat{R} \) is a UFD that we showed in Theorem 36.9. To do so we must analyze division in \( R \) and how it relates to division in \( \hat{R} \). But Lemma 95.26 above has this as a consequence.

**Corollary 95.27.** In the Notation 95.24, suppose that \( a, b \in R \) satisfies \( b \mid a \) in \( \hat{R} \). Then \( b \mid a \) in \( R \).

**Proof.** If \( b \mid a \) in \( \hat{R} \), i.e., \( \hat{R}a \subset \hat{R}b \), then \( \hat{R}a \cap R \subset \hat{R}b \cap R \). Hence \( Ra \subset Rb \) by Lemma 36.9, i.e., \( b \mid a \) in \( \hat{R} \). \( \square \)

We use this to see that if \( x, y \) in our equicharacteristic regular local ring \( R \) are non-associates in \( \hat{R} \), then they are non-associates in \( R \). This is the key to what we need to finish our goal that a equicharacteristic regular local ring is a UFD.

**Proposition 95.28.** In the Notation 95.24, suppose that \( x, y \in R \) have a nonunit factor in \( \hat{R} \). Then \( x \) and \( y \) have a nonunit factor in \( R \).

**Proof.** Since \( \hat{R} \) is a UFD, \( x \) and \( y \) have a greatest common divisor \( \hat{a} \) in \( \hat{R} \) with \( \hat{a} \notin \hat{R}^\times \). Let
\[
x = \hat{a} \hat{a} \quad \text{and} \quad y = \hat{b} \hat{b} \quad \text{with} \quad \hat{a}, \hat{b} \in \hat{R} \quad \text{relatively prime.}
\]
We can approximate
\[
\hat{a} \quad \text{by} \quad \{a_n\} \subset R \quad \text{with} \quad \hat{a} \equiv a_n \pmod{\hat{m}^{n+1}}
\]
\[
\hat{b} \quad \text{by} \quad \{b_n\} \subset R \quad \text{with} \quad \hat{b} \equiv b_n \pmod{\hat{m}^{n+1}}
\]
for every \( n \geq 1 \). So we have
\[
0 = x\hat{b} - y\hat{a} = x(b - b_n + b_n) - y(\hat{a} - a_n + a_n),
\]
hence
\[
(*) \quad xb_n - ya_n = y(\hat{a} - a_n) - x(\hat{b} - b_n).
\]
Let \( \mathfrak{A} = x \mathfrak{m}^{n+1} + y \mathfrak{m}^{n+1} \). The right hand side of \((*)\) lies in \( \hat{R}\mathfrak{A} \cap R = \mathfrak{A} \), using Lemma 95.26. Therefore, we have an equation
\[
xb_n - ya_n = -xr_n + ys_n \quad \text{with} \quad r_n, s_n \in \mathfrak{m}^{n+1}.
\]
Consequently,
\[
x(b_n + r_n) = y(a_n + s_n) \quad \text{in} \quad R.
\]
Therefore we see that in \( \hat{R} \), after cancelling \( \hat{d} \), we have
\[
(**) \quad \hat{a}(b_n + r_n) = \hat{b}(a_n + s_n) \quad \text{in} \quad \hat{R}.
\]
Since \( \hat{a} \) and \( \hat{b} \) are relatively prime in the UFD \( \hat{R} \), we must have \( \hat{a} \) divides every \( a_n + s_n \), i.e.,
\[
(\dagger) \quad a_n + s_n = \hat{a} \hat{u} \quad \text{some} \quad \hat{u} \in \hat{R}.
\]
Reading the equation \((\dagger)\) modulo \( \hat{m}^{n+1} \), we have \( \hat{a} \equiv \hat{a} \hat{u} \mod \hat{m}^{n+1} \), since \( s_n \in \mathfrak{m}^{n+1} \). We cannot have \( \hat{u} \in \hat{m} \), lest \( 1 - \hat{u} \in \hat{R}^\times \). Therefore, \( \hat{u} \in \hat{R}^\times \). Fix \( n >> 0 \) satisfying \( \hat{a} \notin \hat{m}^{n+1} \). Then \( a_n + s_n \) is nonzero in the domain \( \hat{R} \), hence Equation \((\dagger)\) implies that
\[
\hat{a} = \hat{u}^{-1}(a_n + s_n) \quad \text{and} \quad \hat{b} = \hat{u}^{-1}(b_n + r_n).
\]
Therefore,
\[
x = \hat{d} \hat{a} = (a_n + s_n)\hat{u}^{-1} \hat{d} \quad \text{and} \quad y = \hat{d} \hat{b} = (b_n + r_n)\hat{u}^{-1} \hat{d}.
\]
Since \((a_n + s_n) | x \) in \( \hat{R} \), we conclude that \((a_n + s_n) | x \) in \( R \) by Corollary 95.27. It follows by \((**)\) that \( \hat{d} := \hat{u}^{-1} \hat{d} \) in \( R \) and so is a common factor of \( x \) and \( y \) in \( R \). If \( \hat{d} \in R^\times \), then \( \hat{d} = \hat{u} \hat{d} \in \hat{R}^\times \), which is impossible. This proves the proposition. \( \square \)

**Theorem 95.29.** Let \( (R, \mathfrak{m}) \) be an equicharacteristic regular local ring. Then \( R \) is a UFD.

**Proof.** As \( R \) is a Noetherian domain, by Euclid’s Argument (Proposition 29.9), it suffices to show if \( x \in R \) is irreducible, then it is a prime element. So assume that \( x | yz \) with \( y, z \in R \) and \( x \) irreducible in \( R \).

**Case 1.** Suppose that \( x \) and \( y \) have a nonunit common factor in \( \hat{R} \): By Proposition 95.28 above, \( x \) and \( y \) have a nonunit common factor in \( R \). As \( x \) is irreducible in \( R \), we have \( x | y \) in \( R \).
Case 2. Suppose that $x$ and $y$ have no nonunit common factor in $\hat{R}$: Since $\hat{R}$ is a UFD and $x$ and $y$ are relatively prime in $\hat{R}$, we know that $x | z$ in $\hat{R}$. Therefore, $x | z$ in $R$ by Corollary 95.27. □

As a UFD is normal, we have:

**Corollary 95.30.** Let $(R, \mathfrak{m})$ be an equicharacteristic regular local ring. Then $R$ is a normal domain.

**Remark 95.31.** Being a regular local ring is much stronger than just being normal. Let $F$ be a field and $f_1, \ldots, f_m \in F[t_1, \ldots, t_n]$ with $m \geq n$. Suppose that $\tilde{F}$ is an algebraic closure of $F$. Set $\tilde{R} = \tilde{F}[t_1, \ldots, t_n]/(f_1, \ldots, f_m)$. A point $\tilde{x} \in Z_{\tilde{F}}(f_1, \ldots, f_m)$ is called a simple point of $Z_{\tilde{F}}(f_1, \ldots, f_m)$ if the rank of the Jacobian matrix of $\tilde{x}$, i.e., the $m \times n$ matrix $\left( \frac{\partial f_i}{\partial t_j}(\tilde{x}) \right)$, is $n$. Let $\mathfrak{m}_{\tilde{x}}$ be the maximal ideal defined by $\tilde{x}$. Then it is a fact that $\tilde{R}_{\mathfrak{m}_{\tilde{x}}}$ is a regular local ring if and only if $\tilde{x}$ is a simple point. Suppose that $\tilde{F}$ is a perfect field. Then the localization at any maximal ideal $\mathfrak{m}$ of $R = F[t_1, \ldots, t_n]/(f_1, \ldots, f_m)$ is a regular local ring if $\mathfrak{m} = R \cap \mathfrak{m}_{\tilde{x}}$ when $\tilde{x}$ is a simple point in $Z_{\tilde{F}}(f_1, \ldots, f_m)$. In general, this is not true for inseparable extensions.

**Exercises 95.32.**

1. Show if $R$ is a Noetherian ring and $\mathfrak{A}$ is an ideal of $R$, then $G_{\mathfrak{A}}(R)$ is also Noetherian.

2. Show every discrete valuation ring is a regular local ring.

3. Let $R$ be an equicharacteristic local ring with $a, b$ nonzero elements in $R$. Let $\hat{R}$ be as in Notation 95.24. Show that a gcd of $a$ and $b$ in $R$ and in $\hat{R}$ are associates in $\hat{R}$.

4. Let $R$ be a commutative ring and $\mathfrak{A} < R$ an ideal. For each pair of non-negative integers $i \leq j$ let $\theta_{i,j} : R/\mathfrak{A}^j \to R/\mathfrak{A}^i$ be the natural ring epimorphism. Show that there exist a ring $\hat{R}$ and for each $i$ a ring homomorphism $\psi_i : \hat{R} \to R/\mathfrak{A}^i$ such that for all non-negative $i \leq j$
commutes, and if there exist ring homomorphisms $\theta_i : S \to R/\mathfrak{A}^i$ such that for all non-negative $i \leq j$

\[
\begin{array}{ccc}
S & \xrightarrow{\varphi_i} & R/\mathfrak{A}^i \\
\downarrow & & \downarrow \theta_{ij} \\
R/\mathfrak{A}^j & \xleftarrow{\varphi_j} & \end{array}
\]

commutes, then there exists a unique ring homomorphism $\mu : S \to \hat{R}$ such that

\[
\begin{array}{ccc}
S & \xrightarrow{\varphi_i} & \hat{R} \\
\downarrow & & \downarrow \theta_{ij} \\
R/\mathfrak{A}^j & \xleftarrow{\varphi_j} & \end{array}
\]

commutes for all $i \leq j$. Such a $\hat{R}$ is unique up to a unique isomorphism and is called the completion of $R$ relative to the ideal $\mathfrak{A}$.

5. Let $R$ be a Noetherian ring, $\mathfrak{A} < R$ an ideal. Let $\hat{R}$ be the completion of $R$ relative to $\mathfrak{A}$. Show that there exists ring homomorphism $\varphi : R \to \hat{R}$ such that the diagram in 95.33 is valid with $\varphi_i : R \to R/\mathfrak{A}^i$ the canonical surjections and that this map is a monomorphism if $R$ is either a domain or a local ring.

6. In the notation of Exercise 5, show that $\hat{R}$ is the subring of $\prod_{i=0}^{\infty} R/\mathfrak{A}^i$ (a ring with componentwise operations) consisting of all sequences

\[
(\ldots, x_n, \ldots, x_1, x_0 \mid \varphi(x_{n+1} + \mathfrak{A}^{n+1}) = x_n + \mathfrak{A}^n \text{ for all } n \geq 1)
\]

(a subring with component-wise operations).

7. Let $R$ be a Noetherian ring, $\mathfrak{A} < R$ an ideal. Define a topology on $R$, called the $\mathfrak{A}$-adic topology by \{r + \mathfrak{A}^n \mid r \in R, \; n \geq 0\} is a base of open sets for the topology of $R$. Let $\hat{R}$ be the completion of $R$ relative to $\mathfrak{A}$ and $\widehat{\mathfrak{A}} = \hat{R}/\mathfrak{A}$. Show that every Cauchy sequence in the $\widehat{\mathfrak{A}}$-adic topology of $\hat{R}$ converges, i.e., $\hat{R}$ is topologically complete.
8. Let $F$ be a field. Show the completion of $F[t_1, \ldots, t_n]$ is the power series ring $F[[t_1, \ldots, t_n]]$.

9. Let $p$ be a prime in $\mathbb{Z}$, then $\mathbb{Z}$, the completion of $\mathbb{Z}$ relative to the ideal $(p)$ is called the $p$-adic integers [and is historically denoted by $\mathbb{Z}_p$.] Let $\psi_i : \mathbb{Z} \rightarrow \mathbb{Z}/(p^i)$ be the canonical surjections for each positive $i$. Show all of the following:

(i) $\hat{\mathbb{Z}}$ is a domain.

(ii) For each $n > 0$

$$0 \rightarrow \hat{\mathbb{Z}} \xrightarrow{p^n} \hat{\mathbb{Z}} \xrightarrow{\psi_n} \mathbb{Z}/(p^n) \rightarrow 0.$$  

is an exact sequence of additive groups where the map $p^n$ is multiplication by $p^n$. In particular, $\hat{\mathbb{Z}}/p^n\hat{\mathbb{Z}} \cong \mathbb{Z}/(p^n)$ (and are usually identified).

(iii) If $x \in \mathbb{Z}/(p^n)$, then $x$ unit in $\mathbb{Z}/(p^n)$ and only if $p \nmid x$.

(iv) An element $x \in \hat{\mathbb{Z}}$ is a unit if and only if $p \nmid x$.

(v) Every element in $\hat{\mathbb{Z}}$ is of the form $p^r u$ for some $u \in \hat{\mathbb{Z}}^\times$ and $r \geq 0$. In particular, $\hat{\mathbb{Z}}$ is a regular local ring.

(vi) Every element in $\hat{\mathbb{Z}}$ can be uniquely written as a power series $\sum_{i=1}^{\infty} a_i p^i$ with $a_i \in \{0, 1, \ldots, p-1\}$ for all $i$.

96. Addendum: Fibers

One of the first memorable theorems in linear algebra is given vector spaces $V$ and $W$ over a field $F$ with $V$ finite dimensional and a linear transformation $T : V \rightarrow W$, then $\dim V = \dim \text{im} T + \dim \ker T$. As $\ker T = T^{-1}(0)$, the fiber of the zero of $W$ is a subspace of $V$ and $\text{im} T$ a subspace of $W$, this makes sense as a theorem about vector spaces. We wish to generalize the analogous result to maps of irreducible abstract $F$-affine varieties. One problem is that the image of an an affine variety need not be a variety, although its closure is. We write this generalization as a theorem about $F$-affine domains and afterward translate this to the result that we really want.

**Theorem 96.1.** Let $\varphi : A \rightarrow B$ be an injective $F$-algebra homomorphism of $F$-affine domains $A$ and $B$ of dimensions $m$ and $n$ respectively. Then

1. For each $m \in \text{Max}(A)$, each irreducible component of $V(\varphi(m))$ has dimension greater than or equal to $n - m$.

2. There exists an open set $U$ of $\text{Spec}(A)$ such that for all $m$ in $U \cap \text{Max}(A)$, every irreducible component of $V(\varphi(m))$ is of dimension $n - m$. 


We may view $A \subset B$ hence $qf(A) \subset qf(B)$. The transcendence degree of $qf(B)$ over $qf(A)$ is $n - m$. Let $B = F[b_1, \ldots, b_N]$ with $b_1, \ldots, b_{n-m}$ a transcendence basis of $qf(B)$ over $qf(A)$ and $A = F[a_1, \ldots, a_M]$. Hence $m$ of the $a_i$'s are algebraically independent over $F$ and the others polynomials in them (after appropriate clearing of denominators).

Fix $i$, $i = n - m + 1, \ldots, N$. Since $b_i$ is algebraically dependent on $b_1, \ldots, b_{n-m}$ over $qf(A)$, there exist $f_i \in F(a_1, \ldots, a_M, b_1, \ldots, b_{n-m})[t_i]$ satisfying

\[(*) \quad f_i(a_1, \ldots, a_M, b_1, \ldots, b_{n-m}, b_i) = 0\]

We may assume that $f_i$ lies in $F[a_1, \ldots, a_M, b_1, \ldots, b_{n-m}, t_i]$ which we view as $F[a_1, \ldots, a_M][t_i, b_1, \ldots, b_{n-m}]$ by clearing denominators. Let $g_i(a_1, \ldots, a_M) \in F[a_1, \ldots, a_M]$ be the coefficient of $f_i$ of maximal total degree and maximal in the lexicographic order of $t_i, b_1, \ldots, b_{n-m}$ (of non-trivial monomials of $f_i$). Then $V_A(g_i) \subset \text{Spec}(A)$ is closed. Therefore, $V = \bigcup_{i=m-n+1}^N V_A(g_i) < \text{Spec}(A)$ is also closed. In particular, $U = \text{Spec}(A) \setminus V$ is a non-empty open subset. Let $\mathfrak{P} \in U$ and $X = V_B(\mathfrak{p})$ be an irreducible component of $V_B(\mathfrak{P})$. Let $\sim : B \to B/\varphi(\mathfrak{P})$ by $b_i \mapsto \tilde{b}_i$ for $i = 1, \ldots, N$ and $\sim : B/\varphi(\mathfrak{P}) \to B/\mathfrak{p}$ by $\tilde{b}_i \mapsto \tilde{b}_i$ for $i = 1, \ldots, N$. By equation $(*)$, each $\tilde{b}_i$, $i = n - m + 1, \ldots, N$, is algebraically dependent on $\tilde{b}_1, \ldots, \tilde{b}_{n-m}$ over $F$. It follows that $\dim_A V(\mathfrak{P}) = \text{tr deg}_{qf(A)} B/\mathfrak{p} \leq n - m$ and hence by (1), we have equality.

\[\Box\]

**Remark 96.2.** We translate the above theorem in terms of irreducible abstract $F$-affine varieties. If $\varphi : X \to Y$ is a dominant map of irreducible abstract $F$-affine varieties (i.e., the associative map of an $F$-affine monomorphism of $F$-affine domains), then for all closed points $y \in Y$ and irreducible components $Z$ of $\varphi^{-1}(y) := \{ x \in X \mid \varphi(x) = y \}$, we have

$$\dim X \leq \dim Y + \dim Z$$

and there exists an non-empty open set $U \subset X$ such that

$$\dim X = \dim Y + \dim Z$$
for all \( y \in U \). For \( F \)-affine varieties over an algebraically closed field, we get the analogous result with the closed points being elements in \( F^M \) (in the notation of the proof of the theorem) by the Hilbert Nullstellensatz.

**Remark 96.3.** Let \( X \) be a topological space and \( Y \subset X \). We say that \( Y \) is *locally closed* in \( X \) if \( Y \) is open in a closed subset of \( X \). (Equivalently, \( Y \) is the intersection of an open and a closed subset of \( X \).) We say that \( Y \) is *constructible* if \( Y \) is a finite union of locally closed subsets of \( X \). Then it is a fact that if \( \varphi : A \to B \) is an \( F \)-algebra homomorphism of affine \( F \)-algebras, it is always a constructible subset of \( \text{Spec}(A) \) although \( \text{im} \; \varphi \) may not be an affine \( F \)-algebra.

**Exercises 96.4.**

1. Let \( \varphi : A \to B \) be an algebra homomorphism of domains. Using tensor products (cf. Exercises 38.12(25) and (33)), one defines the *scheme theoretic fiber* of \( \varphi : \text{Spec}(B) \to \text{Spec}(A) \) over \( p \in \text{Spec}(A) \) to be \( \text{Spec}(B \otimes_A qf(A/p)) \). Show that \( B \otimes_A qf(A/p) = B_p \otimes_A A/p = B_p/pB_p \).

2. In the previous exercise, show that the map \( (^a\varphi)^{-1}(\mathfrak{p}) \to \text{Spec}(B) \otimes_A qf(A/p) \) by \( \mathfrak{p} \mapsto \mathfrak{p}B_p/p \) is a homeomorphism of topological spaces.

**97. Addendum: Japanese Rings**

Let \( A \) be a Noetherian domain with quotient field \( F \). We say that \( A \) is a *Japanese ring* if whenever \( K/F \) is a finite extension of fields, the integral closure \( A_K \) of \( A \) in \( K \) is a finitely generated \( A \)-module. We saw in Theorem 77.5 that if \( A \) is a Noetherian normal domain with quotient field \( F \) and \( K/F \) a finite separable extension, then the integral closure \( A_K \) of \( A \) in \( K \) is a finitely generated \( A \)-module. In particular, if \( \text{char}(F) = 0 \) a normal domain \( A \) is Japanese. The hypothesis that \( K/F \) is separable is a nontrivial condition if \( \text{char}(F) = p > 0 \). Indeed in the positive characteristic case, \( A \) may not be Japanese, even if \( A \) is normal. In fact, Schmidt-Nagata showed if \( F \) is satisfies \([F : F^p]\) is infinite and

\[
\mathcal{S} := \{ F_\alpha \mid F^p \subset F_\alpha \subset F \text{ are fields with } F_\alpha/F^p \text{ a finite extension} \}
\]

is partially ordered by \( \subset \) with \( \Gamma = \{ \alpha \mid F_\alpha \in \mathcal{S} \} \) and \( A_\alpha = F_\alpha[[t]] \), the formal power series ring with coefficients in \( A_\alpha \). Then \( A_\alpha \) is a DVR with maximal ideal \( A_\alpha \). Further let \( A := \bigcup_{\alpha} A_\alpha \), a ring with the obvious structure. Then it can be shown that \( A \) is a discrete valuation ring but is not Japanese, i.e., the result, in general, is false for dimension one noetherian rings. Of course fields are Japanese. We show even more is
true. To investigate inseparable extensions, we need to use properties of purely inseparable extensions (cf. Exercise 52.10(13)). We begin with the following.

**Lemma 97.1.** Let $A$ be a Noetherian domain with quotient field $F$. Suppose that for all finite, purely inseparable field extensions $E$ of $F$ that $A_E$ is a finitely generated $A$-module. Then $A$ is Japanese.

**Proof.** Let $K/F$ be a finite field extension and $L/K$ a normal closure of $K/F$. As $A_K \subset A_L$ and $A$ is a Noetherian ring, it suffices to show that $A_L$ is a finitely generated $A$-module by Theorem 39.7. So we may assume that $K/F$ is a normal extension. Let $E = K^{G(K/F)}$. Then $E/F$ is Galois, so separable and $E/F$ is purely inseparable by Exercise 52.10(13). As $A_E$ is a finitely generated $A$-module by hypothesis, it is a Noetherian normal domain, so $A_K$ is a finitely generated $A_E$-module by Theorem 77.5. It follows that $A_K$ is a finitely generated $A$-module. □

A Noetherian domain $A$ is called *universally Japanese* if every domain $B$ that is a finitely generated $A$-algebra $B$ is Japanese.

**Theorem 97.2.** Every field is universally Japanese.

**Proof.** Let $A$ be an affine $F$-algebra that is also a domain with quotient field $K$ and $L/K$ a finite field extension. We must show that $A_L$ is a finitely generated $A$-module. By the Noether Normalization Theorem 94.1, there exist $x_1, \ldots, x_n$ in $A$ algebraically independent over $F$ with $F[x_1, \ldots, x_n] \subset A$ integral. As $K/F(x_1, \ldots, x_n)$ is a finite field extension, so is $L/F(x_1, \ldots, x_n)$. Therefore, it suffices to show that $A_L$ is a finitely generated $F[x_1, \ldots, x_n]$-module, i.e., we may assume that $A = F[x_1, \ldots, x_n]$. Since $F[x_1, \ldots, x_n] \cong F[t_1, \ldots, t_n]$, we can assume further that $A = F[t_1, \ldots, t_n]$ and $K = F(t_1, \ldots, t_n)$. Moreover, by Lemma 97.1, we may assume that $L/K$ is purely inseparable. As $A$ is a normal Noetherian domain, we may also assume that $\text{char}(F) = p > 0$.

Since $L/K$ is finite and purely inseparable, there exists a positive integer $m$ satisfying $L^{p^m} \subset K = qf(A)$. Let

$$\mathcal{C} := \text{the set of the } p^m\text{-th roots of all the coefficients of all the } f_i, g_i$$

and
L = K(a_1, \ldots, a_n) \text{ with } a_i = \frac{f_i}{g_i} \text{ with } f_i, g_i \in A, g_i \neq 0 \text{ for all } i
\tilde{F} = F[C] = F[C]
E = \tilde{F}(t_1^{1/p}, \ldots, t_n^{1/p}).

We have $K \subset L \subset E$. As $A$ is a Noetherian ring, it suffices to show that $A_E$ is a finitely generated $A$-module. Let
\[ B = \tilde{F}[t_1^{1/p}, \ldots, t_n^{1/p}] = F[C, t_1^{1/p}, \ldots, t_n^{1/p}]. \]

Each of the finitely many elements, $c \in C, t_1^{1/p}, \ldots, t_n^{1/p}$, is integral over $A$, so $B$ is a finitely generated $A$-module. We also have that $B \subset A_E$ with $qf(B) = E$ and $B = \tilde{F}[t_1^{1/p}, \ldots, t_n^{1/p}]$ a UFD, hence normal. Thus $A_E = B$ and the result follows.

**Corollary 97.3.** Let $F$ be a field and $A$ an affine $F$-algebra that is also a domain with $K = qf(A)$. If $L/F$ be a finite field extension, then $A_L$ is an affine $F$-algebra. In particular, the integral closure of $A$ is a affine $F$-algebra.

A noetherian ring $A$ is called a Nagata ring if $A/p$ is Japanese for every prime ideal $p$ in $\text{Spec}(A)$. It is a deep theorem that a noetherian domain $A$ is a Nagata ring (a property internal to $A$) if and only if it is universally Japanese. In particular, if $A$ is a Dedekind domain of characteristic zero, then $A$ is a Nagata ring as it is Japanese and its quotients determined by maximal ideals are finite fields. Hence any Dedekind domain of characteristic zero is universally Japanese. This is important in arithmetic algebraic geometry.

**98. $C_n$-fields**

We saw that over a finite field, every homogeneous polynomial of degree $d$ in more than $d$ variables has a nontrivial zero. We generalize this result, to classes of fields satisfying every homogeneous polynomial of degree $d$ over the field in more than $d^n$ variables for a fixed $n$ has a nontrivial zero. Examples of fields satisfying this are fields of transcendence degree $n$ over algebraically closed fields. This result will be applicable to the study of non-commutative $F$-algebras to be studied. We do give one example of this in this section, viz., an application to the generalization the Hamiltonian quaternions. Except for using the
Principal Ideal Theorem to begin the first case of an induction proof, this section would have fit in the field theory part of the text.

**Definition 98.1.** Let \( f \in F[t_1, \ldots, t_n] \) be a homogeneous form of degree \( d \). We say that \( f \) is normic if \( f \) has only the trivial zero.

**Example 98.2.** Let \( K/F \) be a finite extensions of fields with \( \mathcal{B} = \{x_1, \ldots, x_n\} \) an \( F \)-basis. The Norm Form of \( N_{K/F} \) is defined to be the homogeneous polynomial \( N_{K/F}(t_1, \ldots, t_n) \in F[t_1, \ldots, t_n] \) of degree \( n \) satisfying \( (a_1, \ldots, a_n) \mapsto N_{K/F}(a_1x_1 + \cdots + a_nx_n) \) for \( a_1, \ldots, a_n \in F \) relative to the basis \( \mathcal{B} \). If \( x = a_1x_1 + \cdots + a_nx_n \) is not zero, then \( N_{K/F}(x) \) is not zero, so \( N_{K/F} \) is normic.

[If \( \text{char } F = p > 0 \) and \( K/F \) is not separable, let \( K/E/F \) with \( E \) the maximal separable extension of \( F \) in \( K \), then \( N_{K/F}(x) \) is defined to be \( (N_{E/F}(x))^{[K:E]} \).]

**Lemma 98.3.** Let \( F \) be a field that is not algebraically closed. Then there exist normic forms over \( F \) of arbitrary large degree.

**Proof.** Let \( K/F \) be of degree \( n > 1 \). By the example, there exists a normic form of degree \( n \). Let \( s > 1 \), and suppose that \( f_s \) is a normic form of degree \( n^s \). Then

\[
\begin{align*}
    f_{s+1}(t_1, \ldots, t_{n^2}) &:= \\
    f_s(f_s(t_1, \ldots, t_n), f_s(f_s(t_{n+1}, \ldots, t_{2n}), \ldots, (f_s(t_{n^2-n+1}, \ldots, t_n))^2))
\end{align*}
\]

is a homogeneous form of degree \( n^{s+1} \) and is a normic form, since any zero of \( f_{s+1} \) would give a zero of \( f_s \). \( \square \)

**Definition 98.4.** A field \( F \) is called a \( C_n \)-field if for every positive integer \( d \) and every homogeneous polynomial \( f \) of degree \( d \) in more than \( d^n \) variables has a non-trivial zero, i.e., we have \( |Z_F(f)| > 1 \). [A \( C_1 \)-field is also called a quasi-algebraically closed field.]

**Example 98.5.** Let \( F \) be a field.

1. \( F \) is algebraically closed field it is clearly a \( C_0 \)-field and conversely no non-algebraically closed field is a \( C_0 \)-field by Example 98.2.
2. If \( F \) is finite, then it is a \( C_1 \)-field by Chevalley-Warning Theorem 65.5.

We give a simple application for \( C_1 \)-fields.

**Proposition 98.6.** Let \( F \) be a \( C_1 \)-field and \( K/F \) a finite field extension. Then \( N_{K/F} : K^\times \to L^\times \) is surjective.
Proof. Suppose that $[K : F] = n$ with $B = \{x_1, \ldots, x_n\}$ an $F$-basis. Let $N_{K/F}(t_1, \ldots, t_n) \in F[t_1, \ldots, t_n]$ denote the norm form based on $B$. Let $0 \neq x \in F$ and $f = N(t_1, \ldots, t_n, t) - xt^n \in F[t_1, \ldots, t_n, t]$, a homogeneous form of degree $n$ in $n + 1$ variables, so has a solution $(a_1, \ldots, a_n, a) \in F^{n+1}$. As the norm form is normic, we cannot have $a = 0$. Let $b_i = a_i/a$ for $i = 1, \ldots, n$. Then $N(b_1x_1 + \cdots + b_nx_n) = N_{K/F}(b_1, \ldots, b_n) = x$. □

The definition of $C_n$-fields arose in Lang’s thesis. One can also define $C_n(d)$-fields, i.e., fields for which all homogeneous polynomials over $F$ of fixed degree $d$ in greater than $d^n$-variables have a non-trivial zero. Let $F$ be a local field, i.e., a field with a complete under a discrete valuation with finite residue class field. This includes all $p$-adic fields arising from completions of number fields under the $p$-adic completion, e.g., the field of $p$-adic numbers $\mathbb{Q}_p$, the quotient field of the $p$-adic integers (cf. Exercise 95.32(9)). Then $F$ is a $C_2(2)$-field. In fact, if $K$ is a number field in which $-1$ is a sum of squares, it is a $C_2(2)$-field (by the theorem of Hasse-Minkowski). Artin, who was Lang’s advisor, conjectured that $p$-adic fields were $C_2$-fields. It was shown that they were $C_2(3)$-fields, but Terjanian showed Artin’s conjecture to be false. However, using model theory, Ax and Kochen, showed that that for a fixed $d$ only finitely many $\mathbb{Q}_p$ were not $C_2(d)$-fields.

In this section, we show that if $F$ is an algebraically closed field, any field of transcendence degree $n$ over $F$ is a $C_n$-field. We apply this theorem to establish a special case of Tsen’s theorem about algebraic extensions of $F(t)$ with $F$ an algebraically closed field. The general case of Tsen’s Theorem will be proven in Section 103.

**Theorem 98.7.** (Lang-Nagata) Let $F$ be a $C_n$-field and $f_1, \ldots, f_r \in F[t_1, \ldots, t_n]$ be homogeneous forms of degree $d$. If $N > rd^n$, then $|Z(f_1, \ldots, f_r)| > 0$.

Proof. If $n = 0$, then $F$ is algebraically closed. By Corollary ?? to the Principal Ideal Theorem, $\dim Z(f_1, \ldots, f_r) \geq N - r > 0$, so the result follows. Therefore, we may assume that $F$ is not algebraically closed, i.e., $F$ is not a $C_0$-field so $n > 0$. In particular, by the lemma, there exists a normic form $\varphi$ of degree $s > r$. Define

$$\varphi_1 := \varphi(f_1(t_1, \ldots, t_N), \ldots, f_r(t_1, \ldots, t_N), f_1(t_{N+1}, t_{2N}), \ldots, f_r(t_{N+1}, t_{2N}), \ldots, f_1(t_{kN+1}, t_{(k+1)N}), \ldots, f_r(t_{kN+1}, t_{(k+1)N}), 0, \ldots, 0)$$
where $rk \leq s < r(k+1)$. Let $[x]$ denote the largest integer in $x \in \mathbb{R}$. So $\varphi_1$ is a homogeneous form in $N\left[\frac{s}{r}\right]$ variables of degree $ds \leq dr\left(\left\lfloor \frac{s}{r} \right\rfloor + 1\right)$ variables. If $n = 1$, i.e., $F$ is a $C_1$-field, then $\varphi_1$ has a nontrivial zero if

\[ (*) \quad N\left[\frac{s}{r}\right] > dr\left(\left\lfloor \frac{s}{r} \right\rfloor + 1\right) \geq ds. \]

Since $N - dr > 0$, equation (*) will hold if $s$ is chosen sufficiently large. As $\varphi$ is normic, a nontrivial zero for $\varphi_1$ gives an element in $Z(f_1 \ldots, f_r)$. If $n > 1$, we can iterate the above process, defining $\varphi_1, \ldots, \varphi_m, \ldots$.

We must show that there exists an integer $m$ satisfying $|Z(\varphi_m)| > 1$. Indeed, suppose such an $m$ exists. Choose a minimal such $m$. As $\varphi_{m-1}$ is normic by the choice of $m$, a non-trivial zero for $\varphi_m$ gives a non-trivial element in $Z(f_1 \ldots, f_r)$.

Consequently, we need only show that there exists a nonzero $m$ satisfying $|Z(\varphi_m)| > 1$. Let

\[ D_m = \deg \varphi_m = d^s = \text{(total degree)} \]

\[ N_m = \text{the number of variables of } \varphi_m = \left[ \frac{N_{m-1}}{r} \right] N. \]

**Claim.** If $m \gg 0$, then $N_m > (D_m)^n$.

If we prove the claim, then we are done by the definition of being a $C_n$-field. We shall, in fact, prove that $N_mD_m^{-n} \to \infty$ as $m \to \infty$.

[If we ignore $\lfloor \cdot \rfloor$, we roughly have $N_m \approx \frac{N_{m-1}}{r} N \approx \frac{N^m}{r^ms}$, so

\[ N_mD_m^{-n} \approx N_mr^{-m}s^{1-m}d^{-mn} \geq r^md^{mn}r^{-m}s^{1-m}d^{-mn} + m + 1 \]

which goes to infinity as $m \to \infty$. Hence the above is plausible.] We have

\[ N_{m+1} = \left[ \frac{N_m}{r} \right] N > \frac{N_m}{r} N > \left( \frac{N_m}{r} - 1 \right) N = N_m \frac{N}{r} - N. \]
Iterating this inequality yields

\[ N_{m+1} > \left( \frac{N_{m-1} N - N}{r} \right) \frac{N}{r} = N_{m-1} \frac{N^2}{r^2} - N \left( 1 + \frac{N}{r} \right) \]

\[ > \cdots > N_1 \left( \frac{N}{r} \right)^m - N \left( 1 + \frac{N}{r} + \frac{N^2}{r^2} + \cdots + \frac{N^{m-1}}{r^{m-1}} \right) \]

\[ = N_1 \left( \frac{N}{r} \right)^m - N \left( \frac{(N/r)^m - 1}{(N/r) - 1} \right) \]

\[ = \left( N_1 - \frac{N_r}{N-r} \right) \left( \frac{N}{r} \right)^m + \frac{N r}{N-r} \]

\[ > \left( N_1 - \frac{N_r}{N-r} \right) \left( \frac{N}{r} \right)^m . \]

As \( N_1 = N[\frac{s}{r}] \), we may choose \( s >> 0 \) to satisfy \( N_1 > N r / (N - r) \).

Hence if we set \( C = N_1 - N r / (N - r) \), we have

\[ \frac{N_{m+1}}{D_{m+1}^n} > \frac{C(N/r)^m}{D_{m+1}^n} = \frac{C(N/r)^m}{d^{m+1} s^n} = \frac{C}{(d s)^n} \left( \frac{N}{r d^m} \right)^m \to \infty \]

as \( m \to \infty \), since \( N > r d^m \).

\[ \square \]

**Theorem 98.8.** (Lang) Let \( F \) be a \( C_n \)-field and \( K/F \) a field extension.

(1) If \( K/F \) is algebraic, then so is \( K \).

(2) \( F(t) \) is a \( C_{n+1} \)-field.

In particular, if \( m = \text{tr.deg}_F K \), then \( K \) is a \( C_{n+m} \)-field.

**Proof.** (1): Let \( f \in K[t_1, \ldots, t_N] \) be a homogeneous form of degree \( d \) and \( N > d^m \). Then there exists a field extension \( K/E/F \) with \( E/F \) finite and \( f \in E[t_1, \ldots, t_N] \). Let \( B = \{x_1, \ldots, x_m\} \), be an \( F \)-basis for \( E \). Let \( t_{ij} \), \( 1 \leq i \leq N \) and \( 1 \leq j \leq m \), be the algebraically independent elements over \( F \) satisfying \( t_i = t_{i1} x_1 + \cdots + t_{im} x_m \) for \( i = 1, \ldots, N \). Substituting these into \( f \) and writing the coefficients of \( F \) in the basis \( B \) yields

\[ f(t_1, \ldots, t_N) = f_1(t_{11}, \ldots, t_{N1}) x_1 + \cdots + f_m(t_{11}, \ldots, t_{N1}) x_m \]

with \( f_i \in F[t_1, \ldots, t_N] \in F(t)[t_{11}, \ldots, t_{N1}] \) homogeneous forms of degree \( d \), \( j = 1, \ldots, m \). Since \( Nm > md^n \), the \( f_i \)'s have a common nontrivial zero by the Lang-Nagata Theorem 98.7. This gives a nontrivial zero for \( f \) over \( K \).

(2): Let \( f \in F(t)[t_1, \ldots, t_N] \) be a homogeneous form of degree \( d \) with \( N > d^{n+1} \). Clearing denominators does not change the condition that \( f \) has or does not have a nontrivial zero, so we may assume that \( f \in F[t][t_1, \ldots, t_N] \). Let \( r = \deg_t f \), the degree of \( f \) in \( t \). Let \( m > 0 \), to
be chosen later. Let $t_{ij}$, $1 \leq i \leq N$, $0 \leq j \leq m$ be the algebraically independent over $F$ satisfying $t_i = t_{i0} + t_{i1}t + \cdots + t_{im}t^m$ for $i = 1, \ldots, N$. Substituting these into $f$ yields

$$f(t_1, \ldots, t_N) = f_0(t_{i0}, \ldots, t_{Nm})t^0 + \cdots + f_m(t_{i0}, \ldots, t_{Nm})t^{dm+r}.$$  

Each $f_i$ is a homogeneous form of degree $d$ in $N(m+1)$ variables. Since $N > d^{n+1}$, we can choose $m >> 0$ to satisfy $N(m+1) > d^n(dm+r+1)$, i.e., $(N - d^{n+1})m > d^n(r+1) - N$. With this choice of $m$, the $f_i$ have a nontrivial common zero that yields one for $f$ over $F(t)$. □

**Corollary 98.9.** Let $F$ be an algebraically closed field and $K/F$ satisfying $\text{tr.deg}_F K = n$. Then $K$ is a $C_n$-field.

We give an application, which we shall generalize later. To do so, we generalize the definition of Hamiltonian quaternions, which in itself, is of mathematical interest.

**Construction 98.10.** We follow the steps in our construction of the Hamiltonian quaternions. Let $F$ be a field of characteristic different than two (leaving some verifications to the reader). ([If the characteristic of $F$ is two, a different generalization holds.] Let $a, b \in F \times$. Let $A$ be a four dimensional vector space over $F$ on basis $\{1, i, j, k\}$ where $1 = 1_F \in V$. One checks that $A$ becomes a (non-commutative) ring by linearly extending

$$(98.11) \quad i^2 = a, \quad j^2 = b, \quad k = ij = -ji$$

(so $k^2 = -ab$) with $0_A = 0_F$ and $1_A = 1_F$, called a generalized quaternion algebra.

$A$ is an $F$-algebra, i.e., $rx = xr$ for all $r \in F, x \in A$ with center $F$ (check). We denote this $F$-algebra by $(a,b)_F$. We also assume when this is written, that it comes with a basis $\{1, i, j, k\}$ satisfying Equation 98.11. So the Hamiltonian quaternions can be written $\mathcal{H} = \left(\frac{-1}{F}, \frac{-1}{F}\right)$.

Check that $\left(\frac{b,a}{F}\right) \cong \left(\frac{a,b}{F}\right)$ as $F$-algebras.

Let $A = \left(\frac{a,b}{F}\right)$. Define $- : A \to A$ by

$$x = x_01 + x_1i + x_2j + x_3k \mapsto \bar{x} = x_01 - x_1i - x_2j - x_3k,$$
with the $x_i \in F$. This map is an anti-isomorphism (just as for the Hamiltonian quaternions) and an involution.

We also have a norm map $N : A \to F$ defined by

$$z \mapsto z\bar{z} = x_0^2 - ax_1^2 - bx_2^2 + abx_3^2$$

be chosen later. Let $t_{ij}$, $1 \leq i \leq N$, $0 \leq j \leq m$ be the algebraically independent over $F$ satisfying $t_i = t_{i0} + t_{i1}t + \cdots + t_{im}t^m$ for $i = 1, \ldots, N$. Substituting these into $f$ yields

$$f(t_1, \ldots, t_N) = f_0(t_{i0}, \ldots, t_{Nm})t^0 + \cdots + f_m(t_{i0}, \ldots, t_{Nm})t^{dm+r}.$$  

Each $f_i$ is a homogeneous form of degree $d$ in $N(m+1)$ variables. Since $N > d^{n+1}$, we can choose $m >> 0$ to satisfy $N(m+1) > d^n(dm+r+1)$, i.e., $(N - d^{n+1})m > d^n(r+1) - N$. With this choice of $m$, the $f_i$ have a nontrivial common zero that yields one for $f$ over $F(t)$. □

**Corollary 98.9.** Let $F$ be an algebraically closed field and $K/F$ satisfying $\text{tr.deg}_F K = n$. Then $K$ is a $C_n$-field.

We give an application, which we shall generalize later. To do so, we generalize the definition of Hamiltonian quaternions, which in itself, is of mathematical interest.

**Construction 98.10.** We follow the steps in our construction of the Hamiltonian quaternions. Let $F$ be a field of characteristic different than two (leaving some verifications to the reader). ([If the characteristic of $F$ is two, a different generalization holds.] Let $a, b \in F \times$. Let $A$ be a four dimensional vector space over $F$ on basis $\{1, i, j, k\}$ where $1 = 1_F \in V$. One checks that $A$ becomes a (non-commutative) ring by linearly extending

$$(98.11) \quad i^2 = a, \quad j^2 = b, \quad k = ij = -ji$$

(so $k^2 = -ab$) with $0_A = 0_F$ and $1_A = 1_F$, called a generalized quaternion algebra.

$A$ is an $F$-algebra, i.e., $rx = xr$ for all $r \in F, x \in A$ with center $F$ (check). We denote this $F$-algebra by $(a,b)_F$. We also assume when this is written, that it comes with a basis $\{1, i, j, k\}$ satisfying Equation 98.11. So the Hamiltonian quaternions can be written $\mathcal{H} = \left(\frac{-1}{F}, \frac{-1}{F}\right)$.

Check that $\left(\frac{b,a}{F}\right) \cong \left(\frac{a,b}{F}\right)$ as $F$-algebras.

Let $A = \left(\frac{a,b}{F}\right)$. Define $- : A \to A$ by

$$x = x_01 + x_1i + x_2j + x_3k \mapsto \bar{x} = x_01 - x_1i - x_2j - x_3k,$$
with the $x_i \in F$. This map is an anti-isomorphism (just as for the Hamiltonian quaternions) and an involution.

We also have a norm map $N : A \to F$ defined by

$$z \mapsto z\bar{z} = x_0^2 - ax_1^2 - bx_2^2 + abx_3^2$$
for all \( z = x_01 + x_1i + x_2j + x_3k \) in \( A \). It satisfies for all \( x, y \in A \):

1. \( 1 = \overline{1} \).
2. \( x + y = \overline{x} + \overline{y} \).
3. \( xy = yx \).
4. \( N(xy) = N(x)N(y) = N(y)N(x) \).

If \( A \) is a generalized quaternion algebra and \( N(x) \neq 0 \) for \( x \in A \), then \( x \) has an inverse, viz., \( x/N(x) \), hence is a division ring if \( N(x) \neq 0 \) for all nonzero \( x \in A \). The difference between the Hamiltonian quaternions and generalized quaternions, is that it is now possible that \( N(x) = 0 \) with \( x \) nonzero which means that \( A \) may not be a division ring.

**Lemma 98.12.** Let \( F \) be a field of characteristic different from two. If \( a, b, x, y \in F^\times \), then

\[
\left( \frac{a, b}{F} \right) \cong \left( \frac{ax^2, by^2}{F} \right).
\]

**Proof.** Let \( \{1, i, j, k\} \) and \( \{1, i', j', k'\} \) be the defining \( F \)-bases for \( \left( \frac{a, b}{F} \right) \cong \left( \frac{ax^2, by^2}{F} \right) \), respectively. As \( (xi)(yj) = xy(ij) = -(yj)(xi) \), we see that the \( F \)-linear map

\[
\varphi : \left( \frac{ax^2, by^2}{F} \right) \to \left( \frac{a, b}{F} \right)
\]

defined by \( i' \mapsto xi, \quad j' \mapsto yj', \) and \( k' \mapsto xyk \) is an \( F \)-linear isomorphism and is checked to be a ring isomorphism (hence an \( F \)-algebra isomorphism).

**Corollary 98.13.** Let \( F \) be a field of characteristic different from two. Then

\[
\left( \frac{1, b}{F} \right) \cong \mathbb{M}_2(F).
\]

**Proof.** If \( \{1, i, j, k\} \) is the \( F \)-basis for \( \left( \frac{1, b}{F} \right) \), then

\[
i \mapsto \begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix} \quad \text{and} \quad j \mapsto \begin{pmatrix} 0 & b \\ 1 & 0 \end{pmatrix}
\]

defines an \( F \)-algebra isomorphism.

**Examples 98.14.** 1. Over the complex numbers, every generalized quaternion \( \mathbb{C} \)-algebra is isomorphic to \( \mathbb{M}_n(\mathbb{C}) \).
2. Over the real numbers, if \( a, b \in \mathbb{R}^\times \),

\[
\left( \frac{a, b}{\mathbb{R}} \right) \cong \begin{cases} \mathcal{H}, & \text{if } a < 0 \text{ and } b < 0 \\ \mathbb{M}_2(\mathbb{R}), & \text{otherwise}. \end{cases}
\]
More generally, we have

**Proposition 98.15.** Let \( F \) be a field of characteristic different from two and \( A = \left( \frac{a, b}{F} \right) \). Then the following are equivalent:

1. \( A \) is isomorphic to \( \mathbb{M}_2(F) \).
2. \( A \) is not a division ring.
3. There exists a nonzero element in \( A \) with zero norm.
4. The element \( b \) is a norm from the field extension \( F(\sqrt{\alpha})/F \).

**Proof.** We need only prove (3) \( \Rightarrow \) (4) and (4) \( \Rightarrow \) (1). If \( a \) is a square, both implications follow, so we may assume not.

(3) \( \Rightarrow \) (4): \( F(\sqrt{\alpha})/F \) is a quadratic extension. If \( x = x_0 1 + x_1 i + x_2 j + x_3 k \) in \( A \) satisfies \( N(x) = 0 \), we have

\[
N_{F(\sqrt{\alpha})/F}(x_2 + \sqrt{\alpha}x_3)b = (x_2^2 - ax_3^2)b = x_0^2 - ax_1^2.
\]

If \( x_2^2 - ax_3^2 \) is zero, this equation implies that \( a \) is a square, so it is not zero. Hence

\[
b = N_{F(\sqrt{\alpha})/F}(x_0 + \sqrt{\alpha}x_1)N_{F(\sqrt{\alpha})/F}(x_2 + \sqrt{\alpha}x_3)^{-1}
\]

and (4) follows.

(4) \( \Rightarrow \) (1): As \( b \) is a norm from \( F(\sqrt{\alpha}) \), so is \( b^{-1} \). Let \( b = y^2 - az^2 \), with \( y, z \in F \). Set \( j' = yj + zij \). Then \( (j')^2 = by^2 - abz^2 = 1 \). Check that \( j'i = -ij' \).

Let \( i' := (1 + a)i + (1 - a)j'i \). Check that \( i'j' = -j'i' \) and \( (i')^2 = (1+a)^2a - (1-a^2)a = 4a^2 \). Then the map \( A \rightarrow \left( \frac{1, 4a^2}{F} \right) \) defined by \( i \mapsto i' \), \( j \rightarrow j' \), \( k \mapsto k' = i'j' \) defines an \( F \)-algebra isomorphism.

Since we know that \( \left( \frac{1, 4a^2}{F} \right) \cong \mathbb{M}_2(F) \), the result follows; \( \square \)

**Corollary 98.16.** Let \( F \) be a field of characteristic not two. A generalized quaternion \( F \)-algebra is either a division \( F \)-algebra or isomorphic to \( \mathbb{M}_2(F) \). In particular, it is an \( F \)-algebra with center \( F \) that is simple as ring and four dimensional as an \( F \)-vector space.

We now give our application to the principal subject of this section. The generalization of this result is called Tsen’s Theorem, which we shall come back to when we study noncommutative algebra.

**Corollary 98.17.** Let \( F \) be a field of transcendence degree one over an algebraically closed field. The every generalized quaternion algebra over \( F \) is isomorphic to \( \mathbb{M}_2(F) \).

**Proof.** Let \( A = \left( \frac{a, b}{F} \right) \) be a generalized quaternion \( F \)-algebra on basis \( \{i, j, k\} \). Define \( N_{A/F} \in F[t_1, \ldots, t_4] \) by \( N_{A/F}(x_0, x_1, x_3, x_4) = \)
$N(x_0 + x_1i + x_2j + x_3k)$. Then $N_{A/F}$ is a homogeneous polynomial of degree four. As $F$ is a $C_1$-field, it has a non-trivial zero. The result follows.

Exercises 98.18.

1. Fill in the details in Construction 98.10.

2. Fill in the details of Proposition 98.15

3. Show that the generalized quaternion algebra $\left( \frac{a, b}{F} \right)$ is isomorphic to $\left( \frac{a, -ab}{F} \right)$. In particular, $\left( \frac{a, a}{F} \right) \cong \left( \frac{a, -1}{F} \right)$.

4. Show if $A = \left( \frac{a, b}{F} \right)$ is a generalized quaternion algebra, then $A$ is not a division $F$-algebra if and only if the quadratic polynomial $at_1^2 + bt_2^2 = 1$ has a non-trivial solution. In particular if $a \neq 0$ or 1, then $\left( \frac{a, 1-a}{F} \right)$ is not a division $F$-algebra.
CHAPTER XVIII

Division and Semisimple Rings

In this chapter, we study the simplest non-commutative rings, viz.,
division rings. More generally, we study matrix rings over division
rings. In the first section, we show that simple left Artinian rings
are up to isomorphism matrix rings over division rings and uniquely
so. We then study products of such rings. This is important as it
is the foundation of the theory of finite group representations that
we shall study in the next chapter. We then investigate another way
of obtaining division rings important in number theory. Finally, we
generalize Wedderburn’s Theorem that finite division rings are fields
to a theorem of Jacobson, showing that a ring is commutative if for
every element \( x \) in the ring, there exists an integer \( n \) satisfying \( x^n = x \).

99. Wedderburn Theory

Modules over a division ring are just vector spaces with a line in
a vector space a ‘simple’ module. As vector spaces are direct sums of
lines, a module over a division ring is a direct sum of ‘simple’ submod-
ules. modules over division rings, i.e., vector spaces are direct sums of
lines. In this section, we shall generalizes this to a module over a simple
left Artinian ring, i.e., any such is a direct sum of ‘simple’ submodules.
We shall also prove Wedderburn’s Theorem that any simple Artinian
ring is, in fact, isomorphic to a matrix ring over a division ring. We
begin with properties of modules.

Definition 99.1. Let \( R \) a nonzero ring. We say
1. A nonzero (left) \( R \)-module is irreducible (or simple) if it contains no
   proper submodules. An irreducible left ideal in \( R \) is also called a
   minimal left ideal.
2. A (left) \( R \)-module is called completely reducible if every submodule
   of it is a direct summand of it.
3. The ring \( R \) is (left) semisimple if it is a completely reducible module
   over itself.

Of course, we have the analogous definitions for right \( R \)-modules. When
we wish to consider right modules we shall write the word right.

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Examples 99.2. 1. Over a division ring, every vector space is completely reducible, since every vector space is free on a basis.

2. Let $D$ be a division ring. Then the $k$th column space of $\mathbb{M}_n(D)$,

$$M_k := \{ A = (a_{ij}) \in \mathbb{M}_n(D) \mid a_{ij} = 0 \text{ if } j \neq k \}.$$ 

is a minimal left ideal of $R$ and is an $n$-dimension vector space over $D$. Moreover, as $\mathbb{M}_n(D)$-modules, we have

$$\mathbb{M}_n(D) = M_1 \oplus \cdots \oplus M_n$$

and $M_i \cong M_j$ for all $i, j$.

Lemma 99.3. Let $\widetilde{M}$ be a nonzero $R$-module and $M_i$ irreducible submodules of $\widetilde{M}$ for $i$ in $I$ (not necessarily non-isomorphic). Let $M$ be the $R$-module $\sum_i M_i$. Then there exists a subset $J$ of $I$ satisfying $M$ is the direct sum $\bigoplus_J M_j$.

Proof. Using Zorn’s Lemma, we see that there exists a subset $J$ of $I$ maximal such that $N = \sum_J M_j = \bigoplus_J M_j$. Let $i_o$ be an element in $I \setminus J$. Then the $R$-module $M_{i_o}$ is irreducible and contains the submodule $N \cap M_{i_o}$, so either $N \cap M_{i_o} = M_{i_o}$, i.e., $M_{i_o} \subset N$ or $N \cap M_{i_o} = 0$, i.e., $\sum_{J \cup \{i_o\}} M_j = \bigoplus_{J \cup \{i_o\}} M_j$. By maximality, this second case does not occur, so $M = N$. \qed

Proposition 99.4. Let $M$ be a nonzero $R$-module. Then the following are equivalent:

1. $M$ is a sum of irreducible $R$-submodules.
2. $M$ is a direct sum of irreducible $R$-modules.
3. $M$ is a completely reducible $R$-module.

Proof. (1) $\Rightarrow$ (2) follows from the lemma.

(2) $\Rightarrow$ (3): Let $N$ be a submodule of $M$ and $M = \bigoplus_I M_i$, with every $M_i$, $i \in I$, irreducible. Using Zorn’s Lemma, there exists a maximal subset $J$ of $I$ such that $N + \bigoplus_J M_j = N \oplus \bigoplus_J M_j$. By the argument in the proof of the lemma, this must be $M$.

(3) $\Rightarrow$ (1): We prove the following:

Claim 99.5. If $M$ is a completely reducible $R$-module and $M_0$ a nonzero submodule of $M$, then there exists an irreducible submodule of $M_0$:

Let $m$ be a nonzero element of $M_0$. Then we have an $R$-epimorphism $\rho_m : R \to Rm$ given by $r \mapsto rm$.

As the the annihilator $\text{ann}_R m$ is a left ideal in $R$, the map $\rho_m$ induces an $R$-isomorphism $\overline{\rho_m} : M/\text{ann}_R m \to Rm$. 

Using Zorn’s Lemma, we know that there is a maximal left ideal \( m \) in \( R \) with \( \text{ann}_R m \subseteq m \). By the Correspondence Principle, \( mm < Rm \) is maximal, hence \( Rm/mm \) is irreducible. By hypothesis, \( M \) is completely reducible, so \( M = M' \oplus mm \) for some submodule \( M' \) of \( M \). Let \( x \) be an element of \( Rm \). Then there exists an \( r \in m \) and an \( m' \in M' \) satisfying 
\[
x = m' + rm \quad \text{so} \quad m' = x - rm \quad \text{lies in} \quad M' \cap Rm.
\]
It follows that 
\[
Rm = (M' \cap Rm) \oplus mm \subset M_0 \quad \text{and} \quad M' \cap Rm \cong Rm/mm \quad \text{is irreducible}.
\]
This establishes the claim. Now let \( M_0 \) be the sum of all the irreducible submodules of \( M \). By the claim \( M_0 \neq 0 \). As \( M \) is completely reducible, \( M = M_0 \oplus M_1 \) for some submodule \( M_1 \). If \( M_1 \neq 0 \), then applying the claim again shows that \( M_1 \) contains an irreducible submodule of \( M \).

This contradicts the choice of \( M_0 \). \(\Box\)

**Corollary 99.6.** Let \( D \) be a division ring. Then \( \mathbb{M}_n(D) \) is semisimple.

**Remark 99.7.** No nonzero ideal in \( \mathbb{Z} \) is irreducible, so \( \mathbb{Z} \) is not semisimple. If \( p > 0 \) is a prime in \( \mathbb{Z} \), then \( \mathbb{Z}/p\mathbb{Z} \) is semisimple as it is a field. More generally, \( \mathbb{Z}/n\mathbb{Z} \) with \( n > 0 \) a product of distinct primes is semisimple. Note that \( \mathbb{Z}/n\mathbb{Z} \), \( n > 0 \) a product of distinct primes, satisfies the descending chain condition, i.e., is an Artinian ring, and has no nonzero nilpotent (left) ideals \( \mathfrak{A} \), i.e., (left) ideals \( \mathfrak{A} \) such that \( \mathfrak{A}^n = 0 \) for some positive integer \( n \). Note also that \( \mathbb{Z}/n\mathbb{Z} \) with \( n > 1 \) not square free is Artinian but has nonzero nilpotent elements and is not semi-simple. (Cf. Exercise 100.7(4) below.)

**Corollary 99.8.** Let \( M_i, \ i \in I, \) be completely reducible \( R \)-modules. Then \( \prod_I M_i \) is completely reducible.

**Corollary 99.9.** Let \( M \) be a completely reducible \( R \)-module, and \( N \) a submodule of \( M \). Then \( N \) and \( M/N \) are completely reducible.

**Proof.** We may assume that \( N \neq 0 \). By Claim 99.5, there exists a submodule \( N_0 \) of \( N \) that is the sum of all the irreducible submodules of \( N \). As \( M \) is completely reducible, \( M = N_0 \oplus N_1 \) for some submodule \( N_1 \). As in the proof of the proposition, we see that \( N = N_0 \oplus (N_1 \cap N) \).

By Claim 99.5, we must have \( N_1 \cap N = 0 \), so \( N = N_0 \) and \( N \) is completely reducible. As \( M = N \oplus N' \) for some submodule \( N' \) of \( M \), \( N' \) is completely reducible by what we have just shown, hence so is \( M/N \cong N' \). \(\Box\)

**Definition 99.10.** Let \( R \) be a ring. An element \( e \) in \( R \) is called an idempotent if \( e^2 = e \). If \( e_1, \ldots, e_n \) are idempotents, they are called orthogonal if \( e_i e_j = \delta_{ij} e_i \) for all \( i \) and \( j \).
Examples 99.11. Let $R$ be a ring.

1. The elements $0$ and $1$ of $R$ are idempotents. They are called the trivial idempotents and are orthogonal.

2. If $e$ is an idempotent so is $1 - e$, and then $e, 1 - e$ are orthogonal idempotents.

3. The sum of any orthogonal idempotents is an idempotent.

4. Let $M = M_1 \oplus M_2$ and $p_i : M \to M$ be the projection $m_1 + m_2 \mapsto m_i$, where $m_i \in M_i$, for $i = 1, 2$. Then $p_1, p_2$ are orthogonal idempotents in the ring $\text{End}_R(M)$.

5. Let $S = \mathcal{M}_n(R)$ and $e_{ij}$ the matrix with $1$ in the $ij$th entry, zero elsewhere. Then $e_{11}, \ldots, e_{nn}$ are orthogonal idempotents in $S$. We also note that $e_{ii}S e_{ii} \cong R$ as rings.

Proposition 99.12. Let $R$ be a nonzero ring. Then the following are equivalent:

1. Every $R$-module is completely reducible.
2. Every short exact sequence in $R$-modules splits.
3. The ring $R$ is semisimple.
4. $R = \bigoplus_{i=1}^n A_i$ for some $n$ and some left ideals $A_i$, with each $A_i$, $i = 1, \ldots, n$, a minimal left ideal; and, furthermore, $A_i = Re_i$, $i = 1, \ldots, n$, with $e_1, \ldots, e_n$ orthogonal idempotents. Moreover if $A$ is a minimal left ideal, then $A = A_i$ for some $i$. (The $A_i$ need not be mutually non-isomorphic.)

Proof. The equivalence of (1) and (2) follows from the corollaries above and (4) $\Rightarrow$ (3) from the proposition. So we need only show (3) $\Rightarrow$ (4).

As $R$ is completely reducible, $R = \bigoplus_j A_j$ for some minimal left ideals $A_j$. There exist nonzero $e_{i_1}, \ldots, e_{i_n}$ in $R$ for some $n$ with $e_{i_j} \in A_{i_j}$ and $1 = e_{i_1} + \cdots + e_{i_n}$. It follows that $R \subset A_{i_1} \oplus \cdots \oplus A_{i_n}$, hence we must have equality. Changing notation, we may assume that $R = A_1 \oplus \cdots \oplus A_n$, with $e_i \in A_i$ for each $i$ and $1 = e_1 + \cdots + e_n$. It follows that $e_j = \sum e_i e_j$ in $A_1 \oplus \cdots \oplus A_n$, hence $e_1, \ldots, e_j$ are orthogonal idempotents. As $0 < Re_j \subset A_j$, we have $Re_j = A_j$, since $A_j$ is irreducible for each $j$. Finally, if $A$ is a minimal left ideal in $R$, then $A \cap A_j > 0$ for some $j$, so $A = A_j$. \hfill $\square$

Let $D$ be a division ring and $R = \mathcal{M}_n(D)$. Then $R$ is a simple ring, i.e., the only 2-sided ideals in $R$ are $0$ and $R$ and $R$ is left and right semisimple, left and right Noetherian (i.e., ACC on left and right ideals), left and right Artinian (i.e., DCC on left and right ideals), since
$R$ is a finite dimensional $D$-vector space. More generally, if $D_1, \ldots, D_r$ are division rings and $A = M_{n_1}(D_1) \times \cdots \times M_{n_r}(D_r)$ then $A$ is left and right semisimple, left and right Noetherian, and left and right Artinian.

One problem that arises when working with vector spaces over non-commutative division rings is that the composition of linear operators and the multiplication of their matrix representations (relative to some fixed bases) of these operators does not correspond, in fact, is reversed. One way around this is to have linear operators and scalars operate on vectors on different sides. Another way is indicated by the following lemma.

**Lemma 99.13.** Let $M$ be an $R$-module. Then the rings $\text{End}_R(\bigsqcup_{i=1}^n M)$ and $M_n(\text{End}_R(M))$ are isomorphic.

**Proof.** Let $N = \bigsqcup_{i=1}^n M_i$ with $M = M_i$ for all $i$. Then we have the usual maps

$$\iota_i : M_i \to N \text{ given by } m \mapsto (0, \ldots, m_i, \ldots, 0)$$

and

$$\pi_j : N \to M_j \text{ given by } (m_1, \ldots, m_n) \mapsto m_j$$

are an $R$-monomorphism and $R$-epimorphism, respectively, satisfying

$$\pi_j \iota_i = \delta_{ij}1_{M_i} \text{ and } \sum_{j=1}^n \iota_j \pi_j = 1_N.$$ 

Let $f \in \text{End}_R(N)$ and $f_{ij} := \pi_i f \iota_j$ in $\text{End}_R(M)$. Define

$$\varphi : \text{End}_R(N) \to M_n(\text{End}_R(M)) \text{ by } f \mapsto (f_{ij}).$$

This map is clearly additive and it is a ring homomorphism, for if $f, g$ lie in $\text{End}_R(N)$, we have

$$(f_{ij})(g_{ik}) = (\pi_i f \iota_j)(\pi_i g \iota_k) = (\pi_i f g \iota_k) \text{ and } 1 \mapsto (\pi_i \iota_j) = I.$$ 

Now define

$$\psi : M_n(\text{End}_R(M)) \to \text{End}_R(N) \text{ by } (f_{ij}) \mapsto \sum_{i=1}^n \sum_{j=1}^n \iota_i f_{ij} \pi_j,$$
an additive map that is a ring homomorphism, for if \((f_{ij}), (g_{ij})\) lie in \(M_n(\text{End}_R(M))\), we have \(\psi(\text{diag}(1_M, \ldots, 1_M)) = 1_N\) and

\[
\psi((f_{il})(g_{kj})) = \psi\left(\sum_{l=1}^{n} f_{il}g_{lj}\right) = \sum_{i=1}^{n} \sum_{j=1}^{n} t_i f_{il}g_{lj} \pi_j
\]

\[
= \sum_{i=1}^{n} \sum_{j=1}^{n} \sum_{l=1}^{n} t_i f_{il} \pi_i g_{lj} \pi_j = \psi(f_{il})\psi(g_{kj}).
\]

Clearly, \(\varphi \circ \psi = 1_{M_n(\text{End}_R(M))}\) and \(\psi \circ \varphi = 1_{\text{End}_R(N)}\), so \(\varphi\) and \(\psi\) are inverse isomorphisms.

For the next lemma, we do write endomorphisms and scalars on different sides.

**Lemma 99.14.** Let \(e\) be an idempotent in the ring \(R\) and define 

\(\rho : eRe \to \text{End}_R(Re)\) by \(eae \mapsto \rho_{eae} : xe \mapsto xe \cdot eae\).

View \(Re\) as a right \((\text{End}_R(Re))\)-module. Then \(\rho\) is an isomorphism of rings.

**Proof.** Check that \(\rho_{eae}\) is an \(R\)-homomorphism of \(Re\) when we write endomorphisms on the right, i.e., \((re)\rho_{eae} = r(e\rho_{eae})\). [Notice when we write it in this way it looks like the associative law.] It is also easily checked that \(\rho\) is a ring homomorphism.

\(\rho\) is a 1 – 1: If \(0 = re \cdot eae = reae\) for all \(r\) in \(R\), then, setting \(r = e\), shows that \(eae = 0\).

\(\rho\) is onto: Let \(f\) lie in \(\text{End}_R(Re)\). Then there exists an element \(a\) in \(R\) satisfying \((e)f = ae\) in \(Re \subset R\). Therefore, for all \(r\) in \(R\), we have

\[
(re)f = (re \cdot e)f = re((e)f) = reae = re \cdot eae = (re)\rho_{eae},
\]

so \(f = \rho_{eae}\). □

The following is essentially an immediate observation, but of great use.

**Lemma 99.15.** (Schur’s Lemma) Let \(M\) be an irreducible \(R\)-module. Then \(\text{End}_R(M)\) is a division ring.

**Proof.** Let \(f\) in \(\text{End}_R(M)\) be nonzero. Then \(\ker f \subsetneq M\) and \(0 < \text{im } f \subsetneq M\). As \(M\) is irreducible, \(\ker f = 0\) and \(\text{im } f = M\), hence \(f\) is an \(R\)-isomorphism. □

We can now classify nonzero simple, left Artinian rings.

**Theorem 99.16.** (Wedderburn’s Theorem) Let \(R\) be a nonzero ring. Then the following are equivalent:

1. \(R\) is a product of simple rings.
2. \(R\) is a division ring.
3. \(R\) is a left Artinian ring.
4. \(R\) is a semisimple Artinian ring.

**Proof.** □
(1) $R$ is simple and left Artinian.
(2) $R$ is simple and semisimple.
(3) There exists a division ring $D$ and a positive integer $n$ such that $R \cong M_n(D)$.

Suppose that $R$ satisfies (3). Then in (3), $D$ is unique up to a ring isomorphism and $n$ is unique. More precisely, if $R$ satisfies (2) and $e$ is a nonzero idempotent in $R$, then the following are true:

(i) All minimal left ideals in $R$ are isomorphic.

(ii) If $\mathfrak{A}$ is a minimal left ideal in $R$ and $D = \text{End}_R(\mathfrak{A})$, then $eRe \cong \mathcal{M}_m(D)$ for some positive integer $m$.

(iii) In (ii), the ring $D$ is a division ring and the center $Z(eRe)$ of $eRe$ is isomorphic to $Z(D)$, the center of $D$.

(iv) If $D$ and $E$ are division rings and $\mathcal{M}_m(D) \cong \mathcal{M}_n(E)$, with $m, n$ positive integers, then $D \cong E$ and $m = n$.

**Proof.** We have already shown that (3) $\Rightarrow$ (1) and (3) $\Rightarrow$ (2).

(1) $\Rightarrow$ (2): As $R$ is left Artinian, there exists minimal left ideal $\mathfrak{A}$ of $R$ by the Minimal Principle. Let

$$0 < \mathfrak{B} := \sum_{r \in R} \mathfrak{A}r \subset R,$$

a (2-sided) ideal. As $R$ is simple, $\sum_{r \in R} \mathfrak{A}r = \mathfrak{B} = R$. Let $\rho_r : \mathfrak{A} \rightarrow \mathfrak{A}r$ be the $R$-epimorphism defined by $a \mapsto ar$. Since $\mathfrak{A}$ is irreducible, either $\rho_r = 0$ or $\rho_r$ is an $R$-isomorphism, i.e., either $\mathfrak{A}r = 0$ or $\mathfrak{A}r \cong \mathfrak{A}$. It follows that $R$ is completely reducible.

(2) $\Rightarrow$ (3): It suffices to show that (2) $\Rightarrow$ (i) — (iv), since 1 is an idempotent and $1R1 \cong R$. The argument to prove (1) $\Rightarrow$ (2) shows:

$$R = \sum_{r \in R} \mathfrak{A}r$$

with $\mathfrak{A}$ a minimal left ideal in $R$ and $\mathfrak{A} \cong \mathfrak{A}r$ for all $r$ in $R$ satisfying $\mathfrak{A}r \neq 0$,

as semisimple rings contain minimal left ideals. We now show conditions (i) — (iv) are satisfied.

(i): As argued before, we see that

There exists a minimal left ideal $\mathfrak{A}$ in $R$.

$R = \mathfrak{A}r_1 \oplus \cdots \oplus \mathfrak{A}r_n$ for some $r_1, \ldots, r_n$ in $R$.

Every minimal left ideal in $R$ is $\mathfrak{A}r_i$ for some $i$, hence is isomorphic to $\mathfrak{A}$.

(ii). As $R$ is completely reducible, there exists a nonzero idempotent $e$ in $R$. As $0 < eRe \subset R$, $eRe$ is also a completely reducible $R$-module.
By Claim 99.5, \( Re \) contains an irreducible submodule, hence a minimal left ideal of \( R \). Using the notation in the proof of \((i)\), we have

\[
Re = \bigoplus_{j=1}^{m} (A_{r_{ij}} \cap Re) \cong \prod_{j=1}^{m} A,
\]

where the \( i_j, 1 \leq j \leq m, \) are those integers satisfying \( A_{r_{ij}} \cap Re \neq 0 \). By Lemmas 99.13 and 99.14 (with \( e = 1 \) as a special case), we have

\[
eRe \cong \text{End}_R(Re) \cong \text{End}_R(\bigoplus_{j=1}^{m} A) \cong \mathbb{M}_{m}(\text{End}_R(A)).
\]

This establishes \((ii)\).

\((iii)\) follows by Schur’s Lemma 99.15 and Exercise 99.20(4).

\((iv)\): Let \( A = \mathbb{M}_m(D), \) \( B = \mathbb{M}_n(E), \) and \( e_{11} = (\delta_{1j}) \) in \( \mathbb{M}_m(D) \). We have \( Ae_{11} \subset \mathbb{M}_m(D) \) is a minimal left ideal. It is unique up to isomorphism by \((i)\). Similarly, if \( e'_{11} = (\delta_{1j}) \) in \( \mathbb{M}_n(E) \), then \( Be'_{11} \) is a minimal left ideal in \( B \), unique up to isomorphism. By Lemma 99.14 and Example 99.11(5), we have

\[
D \cong e_{11}Ae_{11} \cong \text{End}_A(Ae_{11})
\]

\[
E \cong e'_{11}Be'_{11} \cong \text{End}_B(Be'_{11}).
\]

As \( A \cong B \), their unique minimal left ideals (up to isomorphism) must correspond, so \( D \cong E \) by the above. It follows that

\[
m^2 = \dim_D A = \dim_E B = n^2,
\]

so \( m = n \) also. \( \square \)

**Corollary 99.17.** Let \( R \) be a simple ring. then \( R \) is left Artinian if and only if \( R \) is right Artinian if and only if \( R \) is left semisimple if and only if \( R \) is right semisimple.

**Proof.** This is true for \( \mathbb{M}_n(D) \) with \( D \) a division ring. \( \square \)

**Corollary 99.18.** Let \( R \) be a simple, left Artinian ring. Then the following are equivalent:

1. \( R \) is a division ring.
2. The only zero divisor of \( R \) is zero.
3. The only idempotents of \( R \) are 0 and 1.
4. The only nilpotent element in \( R \) is zero.

Let \( F \) be a field. Recall that a nonzero ring \( A \) is called an \( F \)-algebra if there exists a ring homomorphism \( F \to Z(A) \) where \( Z(A) \) is the center of \( A \). This must be a monomorphism and we usually
identify $F$ with $F_1A$. This ring homomorphism makes $A$ into an $F$-module, i.e., a vector space over $F$. We say that the $F$-algebra $A$ is a finite dimensional $F$-algebra if $\dim_F A$ is finite.

**Corollary 99.19.** Let $F$ be a field and $A$ a finite dimensional $F$-algebra. Suppose that $A$ is also a simple ring. Then $A \cong \mathbb{M}_n(D)$ for some division ring $D$ containing $F$ in its center. The division ring $D$ is a finite dimensional $F$-algebra, unique up to isomorphism and $n$ is unique.

**Proof.** The ring $A$ is left Artinian since a finite dimensional vector space over $F$. □

Let $F$ be a field and $A$ a finite dimensional $F$-algebra. If $A$ is also simple and central, i.e., $Z(A) = F$, then $A$ is called a central simple $F$-algebra. By Wedderburn’s Theorem, $A \cong \mathbb{M}_m(D)$ for some division ring $D$. More over $Z(D) = F$ and $D$ is a finite dimensional vector space over $F$. If $D = F$, we say that $A$ is split. If $A$ and $B$ are central simple $F$-algebras, we say that $A$ and $B$ are similar if $A \cong \mathbb{M}_m(D)$ and $B \cong \mathbb{M}_n(E)$ with $D$ and $E$ isomorphic division rings. This is an equivalence relation by Wedderburn’s Theorem. The classes of central simple $F$-algebras under this equivalence relation can be given an abelian group structure and form what is called the Brauer group of $F$, an important group in algebra and number theory. The Brauer group of $\mathbb{R}$ is isomorphic to $\mathbb{Z}/2\mathbb{Z}$ and the Brauer group of $\mathbb{C}$ is trivial, i.e., every central simple algebra over $\mathbb{C}$ splits. This follows from the results in Section 101. A theorem of Tsen shows that the Brauer group of a field of transcendence degree one over $\mathbb{C}$ is also trivial. This is harder to prove. The determination of the Brauer group of a number field is one of the crowning achievements of twentieth century number theory.

**Exercises 99.20.**

1. Let $V$ be a nonzero finite dimensional vector space over a field (or division ring) $F$. Then $V$ is is an $\text{End}_F(V)$-module via evaluation, i.e., $fv := f(cv)$ for all $f$ in $\text{End}_F(V)$ and $v$ in $V$. Show that $V$ is $\text{End}_F(V)$-irreducible.

2. Let $R$ be a commutative ring with ideals $\mathfrak{A}_i$, $i = 1, 2$. Set

$$V_i := \{p \mid p \text{ a prime ideal with } \mathfrak{A}_i \subset p\}.$$  

Suppose that $V_1 \cap V_2 = \emptyset$ and every prime ideal in $R$ lies in $V_1$ or $V_2$, i.e., the set of prime ideals is the disjoint union of $V_1$ and $V_2$. Show that $R$ contains a non-trivial idempotent.
3. Up to isomorphism, show that exists a unique irreducible $R$-module if $R$ is simple and left Artinian.

4. Show if $R$ is a ring, the map $R \to M_n(R)$ given by $r \mapsto rI$ defines an isomorphism between the center $Z(R)$ of $R$ and the center $Z(M_n(R))$ of $M_n(R)$.

100. The Artin-Wedderburn Theorem

In the last section, we classified simple semisimple rings. In this section, we classify semisimple rings. The prototype for the simple case was a matrix ring over a division ring. The prototype for the general case is a finite product of matrix rings over division rings. This generalization is quite useful, as it gives a foundation for that part of the Theory of Group Representations devoted to finite groups over fields of characteristic zero.

If we have a finitely generated completely reducible module over a ring $R$, then we would expect a nice decomposition for it. Indeed, if we look at the case of a finite dimensional vector space then it has a composition series. The analogue gives us the following:

**Lemma 100.1.** Let $M_1, \ldots, M_r$ and $N_1, \ldots, N_s$ be two finite collections of non-isomorphic irreducible $R$-modules. If

$$
M_1^{m_1} \times \cdots \times M_r^{m_r} \cong N_1^{n_1} \times \cdots \times N_s^{n_s}
$$

for some positive integers $m_1, \ldots, m_r, n_1, \ldots, n_s$, then $r = s$ and there exists a permutation $\sigma \in S_r$ such that $M_i \cong N_{\sigma(i)}$ for all $i$. In particular, if $M = N$, then $M_i = N_{\sigma(i)}$ for all $i$.

**Proof.** Let $M = M_1^{m_1} \times \cdots \times M_r^{m_r}$, $N = N_1^{n_1} \times \cdots \times N_s^{n_s}$, and $\varphi : M \to N$ an $R$-isomorphism. Clearly, $\varphi$ and $\varphi^{-1}$ sets up a bijection

$$
\{M_1, \ldots, M_r\} \leftrightarrow \{N_1, \ldots, N_s\},
$$

so $r = s$. Changing notation, we may assume that $M_i \cong N_i$ for all $i$ and we are reduced to showing the following

**Claim.** If $N$ is an irreducible $R$-module and $N^m \cong N^n$, then $m = n$.

Since $D = \text{End}_R(N)$ is a division ring by Schur’s Lemma, and

$$
M_n\left(\text{End}_R(N)\right) \cong \text{End}_R(N^m) \cong \text{End}_R(N^n) \cong M_n\left(\text{End}_R(N)\right)
$$

by Lemma 99.13, we have $m = n$ by Wedderburn’s Theorem 99.16. The last statement is immediate using $\varphi = 1_M$. 

An important generalization of this lemma, the Krull-Schmidt Theorem, occurs when we replace the word irreducible by the word indecomposable, where an $R$-module is called **indecomposable** if it cannot
be written as a non-trivial direct sum of two submodules. We shall not prove this generalization.

**Lemma 100.2.** Let $R$ be a semisimple ring. Then $R$ is left Artinian. More precisely, if $B$ is a non-trivial left ideal of $R$, then there exist unique minimal left ideals $A_1, \ldots, A_m$ in $R$ satisfying $B = A_1 \oplus \cdots \oplus A_m$.

**Proof.** We have shown that $R = A_1 \oplus \cdots \oplus A_n$ for minimal left ideals $A_i$ and if $A$ is a minimal left ideal, then $A = A_i$ for some $i$. It follows, as before, that

$$B = \bigoplus_{i=1}^m (A_{j_i} \cap B)$$

with the $j_i$ satisfying $A_{j_i} \cap B \neq 0$.

The last lemma says that this is unique. As any submodule of $B$ is a direct sum of some of the $A_{j_1}, \ldots, A_{j_m}$, the result follows. □

Our first formulation of the generalization of Wedderburn’s Theorem is the next result. However, in the proof we shall show much more. Since notation will be used in the proof, we shall postpone this more precise statement until after the proof.

**Theorem 100.3.** (Artin-Wedderburn Theorem) Let $R$ be a semisimple ring. Then $R$ is an (internal) direct sum of finitely many simple left Artinian rings, all unique (up to order). In particular, any semisimple right is both left and right Artinian (and left and right semisimple).

**Proof.** We prove this in a number of steps.

**Step 1.** Let $A$ be a minimal left ideal in $R$. Set

$$S_A := \{ B \mid B \subseteq R \text{ a minimal left ideal with } B \cong A \}$$

$$B_A := \sum_{S_A} B.$$

Then the following are true:

(i) $B_A \subseteq R$ is a 2-sided ideal.

(ii) If $A'$ is a minimal left ideal in $R$, then $A \cong A'$ if and only if $B_A B_{A'}$ is not zero.

(i): Let $B \in S_A$ and $x \in R$. As $\rho_x : B \to Bx$ by $y \mapsto yx$ is an $R$-epimorphism, $Bx = 0$ or $Bx \cong B \cong A$. In either case, $Bx \subseteq B_A$, so $B_A$ is a 2-sided ideal.

(ii): ($\Rightarrow$): We know that $A = Re > 0$ with $e$ an idempotent. Then

$$0 \neq e^2 = e \cdot e \in B_A B_A = B_A B_{A'}.$$
(\(\Leftrightarrow\)): Suppose that \(B_3B_{3'}\) is not zero. Then there exist a \(\mathcal{B} \in S_3\) and a nonzero \(b\) in \(\mathcal{B}' \in S_{3'}\) with \(\mathcal{B}b' \neq 0\). As \(\mathcal{B}\) is irreducible, \(\mathcal{B} \rightarrow \mathcal{B}b'\) by \(x \mapsto xb'\) is an \(R\)-isomorphism. As \(0 < \mathcal{B}b' \subset \mathcal{B}'\) with \(\mathcal{B}'\) irreducible, we have
\[
\mathfrak{A} \cong \mathcal{B} \cong \mathcal{B}b' \cong \mathcal{B}' \cong \mathfrak{A}'.
\]

**Step 2.** Let \(R = \mathfrak{A}_1 \oplus \cdots \oplus \mathfrak{A}_q\) with \(\mathfrak{A}_i, i = 1, \ldots, q\), the minimal left ideals in \(R\), arranged such that \(\mathfrak{A}_1, \ldots, \mathfrak{A}_m, m \leq q\), are all the mutually non-isomorphic ones. [So if \(i > m\), there exists a \(j \leq m\) with \(\mathfrak{A}_i \cong \mathfrak{A}_j\) and every minimal left ideal in \(\mathfrak{A}\) is isomorphic to some \(\mathfrak{A}_i\) with \(i \leq m\).]

Set \(B_i = B_{\mathfrak{A}_i}\) for \(1 \leq i \leq m\). Then \(R = \bigoplus_{i=1}^{m} B_i\) and \(B_iB_j = 0\) if \(i \neq j\) (by Step 1). Moreover,

(i) \(R = \bigoplus_{i=1}^{m} B_i\).

(ii) If \(\mathcal{B}\) is a nonzero 2-sided ideal in \(R\), then \(\mathcal{B} = B_i_1 \oplus \cdots \oplus B_i_n\), some \(i_j\), \(1 \leq i_j \leq m\):

(i): Let \(C := B_i \cap \sum_{i \neq j} B_j\). Then
\[
B_iB_j = 0 \text{ so } CB_j = 0 = B_jC \text{ as } C \subset B_i,
\]
\[
B_jB_i = 0 \text{ so } CB_i = 0 = B_iC \text{ as } C \subset \bigoplus_{j \neq i} B_j,
\]
so \(C = RC = 0\).

(ii): As \(R\) is completely reducible, there exists a minimal left ideal \(\mathfrak{A}\) in \(R\) contained in \(\mathcal{B}\), and, in fact, \(\mathcal{B}\) is completely reducible so a sum of minimal left ideals of \(R\). As \(\mathfrak{A}_1, \ldots, \mathfrak{A}_m\) are all the minimal left ideals of \(R\) up to isomorphism, we may assume that \(\mathfrak{A} \cong \mathfrak{A}_i\) lies in \(\mathcal{B}\) and need only show

**Claim.** \(B_i \subset \mathcal{B}\):

Write \(\mathfrak{A} = Re\) with \(e\) a nonzero idempotent. We must show if \(\mathfrak{A}' \in S_3\), then \(\mathfrak{A}' \subset \mathcal{B}\). Let \(a \in \mathfrak{A} = Re = Re \cdot e = \mathfrak{A}e\). Therefore, there exists an element \(b\) in \(\mathfrak{A}\) such that \(a = be = be \cdot e = a \cdot e\), i.e., if \(a \in \mathfrak{A}\) then \(a = ae\). Let \(\varphi : \mathfrak{A} \rightarrow \mathfrak{A}'\) be an \(R\)-isomorphism. Then for every \(a\) in \(\mathfrak{A}\), we have \(\varphi(a) = \varphi(ae) = a\varphi(e)\) in \(\mathfrak{A}'\). Since \(\mathcal{B}\) is a 2-sided ideal, \(\mathfrak{A}' = \mathfrak{A}\varphi(e) \subset \mathcal{B}\), proving the claim.

**Step 3.** Let \(B_i\) be as in Step 2. Then \(B_i\) is a simple, left Artinian ring: We know \(R = B_1 \oplus \cdots \oplus B_m\), so \(1 = f_1 + \cdots + f_m\) for some \(f_i \in B_i\).

**Claim.** The elements \(f_1, \ldots, f_m\) are central orthogonal idempotents, i.e., in addition to the \(f_i\) being orthogonal idempotents, they lie in the center of \(R\):

We know that \(f_1, \ldots, f_m\) are orthogonal idempotents, so we need only show that they are central. Let \(x\) be an element of \(R\). Then we have
x = \sum x_i for some x_i \in B_i for each i. As 1_R x = x 1_R, we must have \( f_i x_i = x_i = x_i f_i \) for all i and \( f_j x_i = 0 = x_i f_j \) for all \( j \neq i \). It follows that \( f_i x = x f_i \) as needed.

It follows by the claim that \( B_i \) is a ring with \( 1_{B_i} = f_i \). Let \( \mathcal{B} \) be a 2-sided ideal in \( B_i \). Since \( B_i B_j = 0 = B_j B_i \) for all \( j \neq i \), we have

\[ \mathcal{B} = B_i \mathcal{B} = (\sum_j B_j) B_i \mathcal{B} = \sum_j B_j \mathcal{B} = R \mathcal{B}, \]

and similarly \( \mathcal{B} = \mathcal{B} R \). Therefore, \( \mathcal{B} \) is a nonzero 2-sided ideal in \( R \).

By Step 2, \( \mathcal{B} \) is a sum of some of the \( B_j \)'s, hence must be \( B_i \), so \( B_i \) is simple. As \( R \) is completely reducible, there exists a minimal left ideal \( \mathfrak{A} \) of \( R \) lying in \( B_i \). If \( \mathfrak{A}' < \mathfrak{A} \) is a left ideal in \( B_i \), then \( B_j \mathfrak{A}' \subset B_j B_i = 0 \) for all \( j \neq i \), so \( \mathfrak{A}' \) is a left ideal in \( R \), hence \( \mathfrak{A}' = 0 \). Therefore, \( \mathfrak{A} \) is a minimal left ideal in \( B_i \). It follows easily that \( B_i \) is semisimple, hence left Artinian.

**Step 4.** Suppose that \( R = \bigoplus_{i=1}^m B_i = \bigoplus_{j=1}^n C_j \) with the \( B_i \)'s as in Step 2 and each \( C_j \) a 2-sided ideal in \( R \) that is also a simple, left Artinian ring. Then \( m = n \) and there exists a permutation \( \sigma \in S_m \) such that \( B_i = C_{\sigma(i)} \) for every \( i \). The \( B_1, \ldots, B_m \) are called the simple components of \( R \) and the decomposition \( R = B_1 \oplus \cdots \oplus B_m \) is called a Wedderburn decomposition of \( R \):

Each \( B_i \cap C_j \) is a 2-sided ideal in both simple rings \( B_i \) and \( C_j \). Therefore, for each \( i \), there exists a \( k \) with \( B_i \cap C_k \) nonzero, hence \( B_i = C_k \).

Similarly, for each \( j \), there exists an \( l \) such that \( B_l \cap C_j \) is nonzero, hence \( B_l = C_j \). The result follows.

**Step 5.** Finish:

Each \( B_i \) is a simple, left Artinian ring, hence also right Artinian (and right semisimple). It follows easily that \( R \) is also right Artinian. \( \square \)

The proof of the theorem also shows the following:

**Theorem 100.4.** (Artin-Wedderburn) Let \( R \) be a nonzero ring. Then the following are equivalent:

1. \( R \) is semisimple.
2. \( R \) is a finite direct product of simple, left Artinian rings.
3. There exist (possibly isomorphic) division rings \( D_1, \ldots, D_m \), some \( m \), unique up to isomorphism (and order), and unique positive integers \( n_1, \ldots, n_m \) and a ring isomorphism

\[ R \cong \mathbb{M}_{n_1}(D_1) \times \cdots \times \mathbb{M}_{n_m}(D_m) \]

up to order.

If \( R \) is a semisimple ring and \( \mathfrak{A}_1, \ldots, \mathfrak{A}_m \) are all the non-isomorphic minimal left ideals in \( R \), then \( \{ \mathfrak{A}_1, \ldots, \mathfrak{A}_m \} \) is called a basic set for \( R \).
We leave the proof of the following two corollaries as exercises:

**Corollary 100.5.** Let \( R \) be a semisimple ring and \( \{A_1, \ldots, A_m\} \) a basic set. If \( M \) is an irreducible \( R \)-module then \( M \cong A_i \) for some \( i \). In particular, if \( R \) is also simple, then, up to isomorphism, there exists a unique irreducible \( R \)-module.

**Corollary 100.6.** Let \( R \) be a semisimple ring and \( \{A_1, \ldots, A_m\} \) a basic sets for \( R \). Let \( B_i = \sum s_{A_i} A \) for \( i = 1, \ldots, r \) and \( M \) a nonzero \( R \)-module. Then \( B_i M \) is a submodule of \( M \) and is a sum of irreducible submodules each isomorphic to \( A_i \). Further, \( M = \bigoplus_{i=1}^m B_i M \).

**Exercises 100.7.**

1. If \( R \) is a semisimple ring and \( M \) a non-trivial \( R \)-module, then \( M \) is a direct sum of irreducible modules, unique up to isomorphism and order.
2. Prove Corollary 100.5.
3. Prove Corollary 100.6.
4. Show that a ring is semisimple if and only if it is left artinian and has no nonzero nilpotent left ideals.

101. Finite Dimensional Real Division Algebras

We have seen that the (Hamiltonian) quaternion algebra \( \mathcal{H} \), the four dimensional real vector space with basis \( \{1, i, j, k\} \), made into a ring by defining a multiplication on this basis by

\[
i^2 = -1, \quad j^2 = -1, \quad ij = -ji = k
\]

and extending this linearly to the whole space defines a division ring. The center \( \{z \in \mathcal{H} \mid yz = zy \text{ for all } y \text{ in } \mathcal{H}\} \) of \( \mathcal{H} \) is \( \mathbb{R} \). Frobenius solved the problem of finding all division algebras \( D \) that are finite dimensional real vector spaces with \( \mathbb{R} \) in its center, i.e.,

\[
\mathbb{R} = \{x \mid xy = yx \text{ for all } y \text{ in } D\}.
\]

We call such algebras **finite dimensional real division algebras**. If \( D \) is such a division algebra, then \( \mathbb{R} \rightarrow D \) given by \( r \mapsto r1_D \) is a ring monomorphism that we view as an inclusion. The result is the following:

**Theorem 101.1.** (Frobenius) Let \( D \) be a finite dimensional real division algebra. Then \( D \) is isomorphic to \( \mathbb{R}, \mathbb{C}, \) or \( \mathcal{H} \).
Proof. We may assume that \( R < D \). As \( D \) is a finite dimensional real vector space, for any element \( x \) in \( D \), we must have \( x \) is a root of a nonzero polynomial in \( \mathbb{R}[t] \) and \( \mathbb{R}[x] \) is a commutative ring. As \( \mathbb{R}[x] \subset D \) with \( D \) a division ring, \( \mathbb{R}[x] \) must be a domain. It follows, just as in the proof of Theorem 47.12, that \( \mathbb{R}[x] = \mathbb{R}(x) \) is a field and a finite extension of \( \mathbb{R} \). (Cf. Exercise 47.24(1).) By the Fundamental Theorem of Algebra, we must have \( \dim_{\mathbb{R}} \mathbb{R}(x) \leq 2 \) and if \( x \) is not a real number, \( \{1, x\} \) is a real basis for \( \mathbb{R}(x) \) and \( \mathbb{R}(x) \cong \mathbb{C} \). In particular, there exists an \( i \in D \) such that \( i^2 = -1 \). Then \( \mathbb{R}(i) \cong \mathbb{C} \), so we can identify \( \mathbb{R}(i) \) and \( \mathbb{C} \), which we do, i.e., we may view \( \mathbb{R} \subset \mathbb{C} \subset D \). Therefore, \( D \) is a (left) finite dimensional complex vector space. Let \( T : D \to D \) be the \( \mathbb{C} \)-linear transformation given by \( x \mapsto xi \). Then \( T \) satisfies \( T^2 = -1_D \), so the minimal polynomial \( q_T \) of \( T \) satisfies \( q_T \mid t^2 + 1 \) in \( \mathbb{C}[t] \). It follows that the only possible eigenvalues of \( T \) are \( \pm i \) and at least one of them must be an eigenvalue. Let \( E_T(\alpha) = \{v \in D \mid Tv = \alpha v\} \) for \( \alpha \in \mathbb{C} \), so either \( E_T(i) \neq 0 \) or \( E_T(-i) \neq 0 \). Since \( E_T(i) \cap E_T(-i) = 0 \), we have \( E_T(i) \oplus E_T(-i) \subset D \). Let \( x \) be an element in \( D \). Then we have

\[
\begin{align*}
T(x - ixi) &= xi + ix = i(x - ixi), \text{ so } x - ixi \text{ lies in } E_T(i) \\
T(x + ixi) &= xi - ix = -i(x - ixi), \text{ so } x + ixi \text{ lies in } E_T(-i).
\end{align*}
\]

As

\[x = \frac{1}{2}(x - ixi) + \frac{1}{2}(x + ixi) \text{ lies in } E_T(i) + E_T(-i), \]

we have \( D = E_T(i) \oplus E_T(-i) \) as a complex vector space.

If \( x \in E_T(i) = \{x \in D \mid ix = xi\} \), then \( \mathbb{C}[x] \) is a commutative division ring, i.e., a field. Therefore, by the Fundamental Theorem of Algebra, \( E_T(i) = \mathbb{C} \), so \( \dim_{\mathbb{C}} E_T(i) = 1 \). Consequently, if \( E_T(-i) = 0 \), we must have \( D = \mathbb{C} \). So we may assume that \( E_T(-i) \) is nonzero. If \( x \) and \( y \) lie in \( E_T(-i) \), then

\((^*)\)

\[xyi = x(-i)y = ixy,\]

so \( xy \) lies in \( E_T(-i) \). Let \( y \) be a nonzero element in \( E_T(-i) \). Define a \( \mathbb{C} \)-linear transformation

\[\rho_y : D \to D \text{ by } x \mapsto xy.\]

As \( D \) is a division ring, \( \rho_y \) must be a monomorphism (why?). By \((^*)\), we have

\[\rho_y|_{E_T(-i)} : E_T(-i) \to E_T(i).\]

As the linear transformation \( \rho_y|_{E_T(-i)} \) is injective and \( \dim_{\mathbb{C}} E_T(i) = 1 \), we must also have \( \dim_{\mathbb{C}} E_T(-i) = 1 \).

Claim: If \( x \) is a nonzero vector in \( E_T(-i) \), then \( x^2 \) is real and \( x^2 < 0 \):
We know that $x^2 \in \mathbb{R}(x) \cap E_T(i) = \mathbb{R}(x) \cap \mathbb{C}$ by (*). As $\mathbb{R}[x]$ is a two dimensional real vector space on basis \{1, x\} with $x$ not lying in $\mathbb{C}$, if $a + bx$, $a, b$ real, lies in $\mathbb{C}$, we must have $b = 0$. It follows that $x^2$ is real. If $x^2 > 0$, then $x^2 = \theta^2$ for some real number $\theta$. This means that $\pm \theta, x$ would be three distinct roots of $t^2 - x^2$ in $\mathbb{R}[t]$, which is impossible. Therefore, $x^2 < 0$ and the claim is established.

Now let $j = x/\sqrt{|x^2|}$. Then $j^2 = -1$. Set $k = ij$, Then \{j, k\} is a real basis for the vector space $E_T(-i)$ — check — and \{1, i, j, k\} is a basis for the real vector space $D$ satisfying $i^2 = -1 = -j^2$ and $k = ij = -ji$. If follows that $D \cong \mathcal{H}$. □

102. Cyclic Algebras

Wedderburn’s Theorem shows that any central simple algebra is a matrix ring over a division ring. However, it does not give any indication on how to find division rings. Wedderburn was also interested in discovering new division rings. The simplest construction of new central simple algebras became important in Number Theory. In this section, we introduce this construction. Behind its foundation is Hilbert’s Theorem 90 (and its cohomological formulation that we shall not discuss).

We begin with the construction of new rings.

**Definition 102.1.** Let $R$ be any nonzero ring and $\sigma$ a ring automorphism of $R$. Define the twisted polynomial ring $R[t, \sigma]$ to be the ring with the usual addition of polynomials with multiplication given by

$$\left( \sum_i a_i t^i \right) \left( \sum_j b_j t^j \right) := \sum_{i,j} a_i \sigma^i(b_j)t^{i+j}$$

where the $a_i, b_j$ lie in $R$. So, in general, $t$ does not commute with elements in $R$, rather we have $tb = \sigma(b)t$ for all $b$ in $R$.

**Remarks 102.2.** Let $R$ be a nonzero ring and $\sigma$ a ring automorphism of $R$.

1. If for any $a$ and $b$ in $R$, satisfying $ab = 0$, we have $a = 0$ or $b = 0$, i.e., $R$ is a non-commutative domain, then so is the twisted polynomial ring $R[t, \sigma]$.

2. Let $R = K$ be a field. Then the left division algorithm holds in $K[t, \sigma]$. In particular, every left ideal is principal (i.e., $K[t, \sigma]$ is a left PID). If $F \subset K(\sigma)$ is a subfield and $f$ is a polynomial in $F[t]$, then $K[t, \sigma]f$ is a two sided ideal, so $K[t, \sigma]/K[t, \sigma]f$ is a ring.
Definition 102.3. Let $K/F$ be a (finite) cyclic field extension of degree $n$, i.e., Galois with cyclic Galois group, with $G(K/F) = \langle \sigma \rangle$ and $f = t^n - a$ a polynomial in $F[t]$ with $a$ nonzero. Then $(K/F, \sigma, a) := K[t, \sigma]/K[t, \sigma]f$ is called a cyclic algebra of degree $n$ over $F$. Equivalently, using the canonical epimorphism $\bar{-} : K[t, \sigma] \to K[t, \sigma]/K[t, \sigma]f$, we have $(K/F, \sigma, a)$ is a vector space over $K$ on basis $\{1, x, \ldots, x^{n-1}\}$ with $x = \bar{t}$ and satisfying

$$x^n = a \quad \text{and} \quad x^\beta = \sigma(\beta)x \quad \text{for all } \beta \in K.$$ 

It is convenient to use both of these formulations, so we switch between them with this notation without further comment.

Note that a cyclic algebra $(K/F, \sigma, a)$ of degree $n$ over $F$ is a finite dimensional $F$-algebra of $F$-dimension $n^2$.

Examples 102.4. 1. The Hamiltonian quaternions $\mathcal{H}$ is the cyclic algebra $(\mathbb{C}/\mathbb{R}, -1, -1)$ over $\mathbb{R}$ of degree 2. Indeed, viewing $\mathbb{C} = \mathbb{R}(\sqrt{-1})$, then $j$ is the unique element in $\mathcal{H}$ satisfying $j\alpha j^{-1} = \overline{\alpha}$ for all, i.e., $j\alpha = \overline{\sigma j}$ for all $\alpha \in \mathbb{C}$ (as it is true for $i$) and $j^2 = -1$.

2. More generally, the generalized quaternions in Construction 98.10 are cyclic algebras over a field $F$ of characteristic different from two (and in characteristic two, a different construction produces an $F$-algebra of generalized quaternions which we have not constructed).

Indeed, let $A = \left( \frac{a, b}{F} \right)$ be a generalized quaternion algebra with $a$ not a square in $F$ with $F$-basis $\{1, i, j, k\}$. Let $K = F(\sqrt{a}) = F(i)$ and $G(K/F) = \langle \sigma \rangle$. Then $j \in A^\times$ satisfies $j\alpha j^{-1} = \alpha$ for all $\alpha \in K$, i.e., $j\alpha = \overline{\sigma j}$, just as above, with $j^2 = b$.

3. We now generalize the construction of generalized quaternions. Let $F$ be a field of characteristic zero or of prime degree not dividing $n$ and containing a primitive $n$th root of unity. Suppose that $K/F$ is a cyclic extension of degree $n$ with $G(K/F) = \langle \sigma \rangle$. Then $K = F(\alpha)$ with irreducible $m_F(\alpha) = t^n - a$ in $F[t]$, so $\alpha^n = a$ and $\sigma(\alpha) = \zeta \alpha$ for some primitive $n$th root of unity $\zeta$ by Theorem 59.20 and its proof. Let $A$ be an $F$-algebra of dimension $n^2$ on basis $\{\alpha^i \beta^j \mid 0 \leq i, j \leq n - 1\}$ with $K \subset A$ satisfying

$$A = \prod_{i,j=0}^{n-1} F\alpha^i \beta^j$$

$$yz = \zeta zy$$

$$\alpha^n = a, \quad \beta^n = b$$
Then $\beta \theta \beta^{-1} = \sigma(\theta)$, i.e., $\beta \theta = \sigma(\theta) \beta$, for all $\theta \in K$, so $A$ is a cyclic $F$-algebra. Conversely, if $A$ is a cyclic $F$-algebra with $K \subset A$, then we shall see in Section 103 that the Skolem-Noether Theorem 103.18 will imply that there exists $\beta \in A$ satisfying $\beta^{-1} \theta \beta^{-1} = \sigma(\theta)$, i.e., $\beta \theta = \sigma(\theta) \beta$, for all $\theta \in K$. In particular, $\beta^n = b$, for some $b \in F$, and $A = (K/F, \sigma, a)$. This cyclic algebra is usually written $(\frac{a}{b}, F, \zeta)$.

We shall also see that if $F$ has a cyclic extension of degree $n$, then $F_n(F)$ is also an example of a cyclic algebra. To do this, we show that cyclic algebras are central simple algebras.

**Proposition 102.5.** Let $R = (K/F, \sigma, a)$ be a cyclic algebra over $F$ of degree $n$. Then $R$ is a central simple algebra over $F$. Moreover, the centralizer $Z_R(K) := \{ x \in R \mid xb = bx \text{ for all } b \in K \}$ is $K$ and the only field $L$ satisfying $K \subset L \subset R$ is $K$.

**Proof.** We first show that $R$ is simple. Write

$$R = K 1_R \oplus \cdots \oplus K x^{n-1}$$

$$x^n = a \text{ with } a \in F \text{ and } x\alpha = \sigma(\alpha)x \text{ for all } \alpha \in K.$$

Let $0 < \mathfrak{A} \subset R$ be a 2-sided ideal. Choose a nonzero element $z \in \mathfrak{A}$ satisfying

$$z = \alpha_1 x^{i_1} + \cdots + \alpha_r x^{i_r} \text{ with } \alpha_1, \ldots, \alpha_r \in K^\times, 0 \leq i_1 < \cdots < i_r < n$$

and $r$ minimal. If $r = 1$, then $z$ is a unit in $R$ as $\alpha_i x^{i_i}$ are units, so we may assume that $r > 1$. As $\sigma^{i_1} \neq \sigma^{i_r}$, there exists and element $\beta$ in $K$ satisfying $\sigma^{i_1}(\beta) \neq \sigma^{i_r}(\beta)$. In the ideal $\mathfrak{A}$, we have

$$\begin{align*}
(1) \quad z\beta &= (\alpha_1 x^{i_1} + \cdots + \alpha_r x^{i_r})\beta \\
&= \alpha_1 \sigma^{i_1}(\beta)x^{i_1} + \cdots + \alpha_r \sigma^{i_r}(\beta)x^{i_r} \\
(2) \quad \sigma^{i_1}(\beta)z &= \alpha_1 \sigma^{i_1}(\beta)x^{i_1} + \cdots + \alpha_r \sigma^{i_1}(\beta)x^{i_r}.
\end{align*}$$

Subtracting (2) from (1) yields the nonzero element

$$z\beta - \sigma(\beta)z = \alpha_2 (\sigma^{i_2}(\beta) - \sigma^{i_1}(\beta))x^{i_2} + \cdots + \alpha_r (\sigma^{i_r}(\beta) - \sigma^{i_1}(\beta))x^{i_r},$$

contradicting the minimality of $r$. Therefore, $R$ is simple.

Suppose that $z = \sum_{i=0}^{n-1} \alpha_i x^i$ is nonzero in $R$ and commutes with all $\beta$ in $K$. Then $\alpha_i \sigma(\beta)x^i = \alpha_i \beta x^i$ for $i = 0, \ldots, n-1$. If $\beta$ does not lie in $F$, then $\alpha_i = 0$ for $i = 1, \ldots, n-1$ and $\beta$ must lie in $K$. If $K \subset L \subset R$ is a subfield, then $L \subset Z_R(K) = K$.

Finally, if $z \in Z(R)$, then $z \in Z_R(K) = K$. As $z\beta = \sigma(\beta)z$ for all $\beta \in K$, we have $z \in K^{(\sigma)} = F$. $\square$
A field $L$ lying in the cyclic algebra $R = (K/F, \sigma, a)$ not properly contained in any larger in $R$ is called a maximal subfield of $R$. So $K$ is one such. There can be many non-isomorphic ones.

The proposition allows us to show:

**Computation 102.6.** Let $K/F$ be a cyclic extension of fields of degree $n$ with Galois group $\langle \sigma \rangle$ and $R = (K/F, \sigma, 1)$. We show that $R \cong \mathbb{M}_n(F)$:

Since $t^n - 1 = (t-1)(t^{n-1} + \cdots + t + 1)$ in $K[t, \sigma]$ and $R = K[t, \sigma]/(t^n - 1)$, we have $t^n - 1$ lies in the principal ideal $R(t-1)$, a maximal left ideal in $R$ by Remark 102.2(2). It follows that $M := K[t, \sigma]/K[t, \sigma](t-1)$ is a simple $R$-module of $F$-dimension $n$. Since $R$ is a simple ring, the $F$-algebra homomorphism $\rho : R \to \text{End}_F(M)$ defined by $r \mapsto \rho_r : m \mapsto mr$ must be a monomorphism. Since $\dim_F R = \dim_F \text{End}_F(M)$ and $\mathbb{M}_n(F) \cong \text{End}_F(M)$, the result is established.

**Corollary 102.7.** Let $R = (K/F, \sigma, a)$ be a cyclic algebra over $F$ of prime degree $p$. Then $R$ is either a division ring or $R \cong \mathbb{M}_p(F)$.

Recall that a central simple $F$ algebra is called split if it is isomorphic to a matrix ring over $F$.

**Proof.** By Wedderburn's Theorem 99.16, $R \cong \mathbb{M}_r(D)$ for some division ring $D$, so $p^2 = r^2 \dim_F(D)$. As $p$ is a prime number, $p = r$ and $\dim_F(D) = 1$ splits, (i.e., $R$ splits) or $R$ is a division algebra (i.e., $r = 1$ and $R = D$).

**Theorem 102.8.** Let $R = (K/F, \sigma, a)$ be a cyclic algebra of degree $n$ over $F$. Then $R \cong \mathbb{M}_n(F)$ if and only if $a \in N_{K/F}(K^\times)$.

**Proof.** ($\Leftarrow$): Suppose that $a \in N_{K/F}(K^\times)$. Then there exists an $\alpha \in K^\times$ such that $aN_{K/F}(\alpha) = 1$. Let $y = ax$ where $x$ is the image of $t$ in $K[t, \sigma]$ under the canonical epimorphism. Set $y = \alpha x$ in $R$. As $x\alpha = \sigma(\alpha)x$ and $x^n = a$, we have

$$y^n = (ax)^n = \sigma^{n-1}(\alpha) \cdots \sigma(\alpha)ax^n = N_{K/F}(\alpha)a = 1,$$

so, if $\beta \in K$,

$$y\beta = ax\beta = \alpha\sigma(\beta)x = \sigma(\beta)x = \sigma(\beta)y.$$

The $F$-algebra map $R \to (K/F, \sigma, 1)$ sending $x$ to $y$ and fixing $K$ is therefore an isomorphism. It follows by the computation that $R \cong \mathbb{M}_n(F)$.

($\Rightarrow$): Suppose that $R \cong \mathbb{M}_n(F)$. Then $R$ has a simple left $R$-module $M$ of dimension $n$ over $F$. As $K[t, \sigma]$ is a left PID, we have $M \cong \mathbb{M}_n(F)$. 

\[ K[t, \sigma]/K[t, \sigma]f \] for some \( f \in K[t, \sigma] \) with \( K[t, \sigma]f \supseteq K[t, \sigma](t^n - a) \).

Since
\[ n = \dim_K M = (\dim_F K)(\deg f), \]
we must have \( \deg f = 1 \), i.e., \( f = t - c \) for some \( c \) in \( K \). As \( K[t, \sigma]f \supset K[t, \sigma](t^n - a) \),
\[ t^n - a = (b_{n-1}t^{n-1} + \cdots + b_1t + b_0)(t - c), \]
for some \( b_0, \ldots, b_{n-1} \) in \( K \). Multiplying out and comparing coefficients yields
\[ b_{n-1} = 1, \quad b_{n-2} = \sigma^{n-1}(c), \quad b_{n-3} = \sigma^{n-1}(c)\sigma^{n-2}(c), \quad \vdots \]
\[ b_0 = \sigma^{n-1}(c)\cdots\sigma(c). \]
So \( a = b_0c = \sigma^{n-1}(c)\cdots\sigma(c)c = N_{K/F}(c) \) as needed. [Cf. this proof with that of Hilbert Theorem 90.]

Let \( K/F \) be a cyclic Galois field extension of degree \( n \) with Galois group \( \langle \sigma \rangle \) and \( \sigma : F^\times \to F^\times/N_{K/F}(F^\times) \) be the canonical map. Wedderburn showed that \( (K/F, \sigma, a) \) is a division ring if \( a \) has order \( n \) in \( F^\times/N_{K/F}(K^\times) \). This generalizes Corollary 102.7. A deep theorem in Algebraic Number Theory shows that every division algebra over a number field is a cyclic algebra. This is not true in general.

There is a more general construction for an arbitrary finite Galois extension \( K/F \) called a crossed product algebra that also produces central simple algebras that we shall study in §104. It is constructed as follows. Let \( K/F \) be a finite Galois extension of fields. Let \( R \) be a vector space over \( K \) of dimension \( [K : F] \) on basis \( B := \{u_\sigma \mid \sigma \in G(K/F)\} \).

We define multiplication on \( R \) as follows: Let \( \sigma, \tau, \eta \) be arbitrary elements in \( G(K/F) \) and set
\[ u_\sigma \beta = \sigma(\beta)u_\sigma \text{ for all } \beta \in K \]
\[ u_\sigma u_\tau = \alpha_{\sigma,\tau}u_{\sigma\tau} \text{ for some } \alpha_{\sigma,\tau} \in K \]
\[ \alpha_{\sigma,\tau}\alpha_{\sigma\tau,\eta} = \sigma(\alpha_{\tau,\eta})\alpha_{\sigma,\tau\eta}. \]
The last relation is to guarantee associativity. It also shows that \( \alpha_{1,\sigma} = \alpha_{1,1} \) and \( \alpha_{\sigma,1} = \sigma(\alpha_{1,1}) \) so the one of \( R \) is \( \alpha_{1,1}^{-1}u_1 \). One also checks that \( R \) is an \( F \)-algebra, i.e., \( a(\alpha\beta) = (aa\alpha)\beta = \alpha(a\beta) \) for all \( a \in F \) and \( \alpha, \beta \in R \). It is a fact that, although every equivalence class of central
simple algebras contains a crossed product algebra, not every central simple algebra is one. So the theory is still not finished.

**Exercises 102.10.**

1. Let $R$ be a ring and $\sigma$ a ring automorphism of $R$. Define the **twisted power series ring** $R[[t, \sigma]]$ over $R$, to be the ring with the usual addition of (formal) power series and multiplication induced by $tx = \sigma(x)t$ for all $x$ in $F$. Show if $R$ is a non-commutative domain, so is $R[[t, \sigma]]$.

2. Let $R$ be a field and $\sigma$ a ring automorphism of $R$. Define the **twisted Laurent series ring** $R((t, \sigma)) := \{ \sum_{-\infty}^{\infty} a_i t^i \mid a_i \in R \text{ with } a_i = 0 \text{ for almost all negative } i \}$ over $R$, to be the ring with the usual addition of (formal) Laurent series (i.e., $\sum_{-\infty}^{\infty} a_i t^i + \sum_{-\infty}^{\infty} b_i t^i = \sum_{-\infty}^{\infty} a_i b_i t^i$) and multiplication induced by $tx = \sigma(x)t$ for all $x$ in $R$. If $R$ is a division ring, show that $R((t, \sigma))$ is also a division ring and never a field if $\sigma$ is not the identity.

3. Let $R$ be a field and $\sigma$ a ring automorphism of $R$ of finite order $n$. Show that the center of $R((t, \sigma))$ is $R((t^n)) = R((t^n, 1_R))$. In particular, $R((t, \sigma))$ is not finite dimensional over its center.

4. Let $F$ be a field of characteristic zero of positive degree not dividing $n$ containing a primitive $n^{\text{th}}$ root of unity. Suppose that $K/F$ is a cyclic extension of degree $n$ and $t^n - a$ and $t^n - b$ are irreducible polynomials in $F[t]$. Show that $\left( \frac{b, a}{F, \zeta^{-1}} \right) \cong \left( \frac{a, b}{F, \zeta} \right)$.

5. Show equations 104.7 define a ring.

6. Show the ring defined by equations 104.7 is a central simple $F$-algebra.

### 103. Central Simple Algebras

In this section we shall study simple $F$-algebras more carefully, with $F$ a field, under the restriction that the center $Z(A) := \{ x \in A \mid yx = xy \text{ for all } y \in A \} \subset A$ of $A$ is $F$ and, furthermore, that $A$ is a finite dimensional $F$-algebra (as an $F$-vector space) that we studied in §90. As such an $A$ is left Artinian over $F$, Wedderburn’s Theorem 99.16 applies, so $A \cong \mathcal{M}_n(D)$ for some unique integer $n$ and division $F$-algebra $D$, unique up to isomorphism. We use this to study the collection of such algebras and see that it gives an appropriate generalization of algebraic field theory.
In this section we shall need to use the definition and properties of tensor products of modules and algebras. Some of these were given as exercises in Exercises 38.12(25)–(33). A full discussion can be found later in the book in Section 116 below. In particular, let \( A \) be an \( R \)-algebra with \( R \) a commutative ring (every ring is an algebra over some commutative ring) and \( M, N \) (left) \( A \)-modules. Then the tensor product \( M \otimes_R N \) is a (left) \( R \)-module generated by \{ \( x \otimes y \mid x \in M, y \in N \) \} with the \( R \)-action induced by \( r(x \otimes y) = rx \otimes y = x \otimes ry \) for all \( x \in M, y \in N \), and \( r \in A \). If \( M \) is a left \( A \)-module, then \( A \otimes_R M \) is canonically isomorphic to \( M \) as a left \( A \)-module and \( M \otimes_R N \) is an \( A \)-module. In particular, if \( B \) is another \( R \)-algebra, then the \( R \)-module \( B \otimes_R M \) becomes a \( B \)-module by \( b(b_1 \otimes m) = bb_1 \otimes m \) for all \( b, b_1 \in B \), \( m \in M \). In addition, \( B \otimes_R A \) is an \( R \)-algebra with multiplication induced by \( (b_1 \otimes a_1)(b_2 \otimes a_2) = b_1b_2 \otimes a_1a_2 \) for all \( a_1, a_2 \in A \), \( b_1, b_2 \in B \). For example, if \( K/F \) is an extension of fields an \( F \)-vector space \( V \) induces a \( K \)-vector space \( K \otimes_F V \) (extension of scalars).

Let \( A \) be an \( R \)-algebra. Define the \textit{opposite algebra} \( A^{\text{op}} \) of \( A \) as follows: Let \( A^{\text{op}} := \{a^{\text{op}} \mid a \in A \} \) and \( (\ )^{\text{op}} : A \rightarrow A^{\text{op}} \) be the bijection given by \( a \mapsto a^{\text{op}} \). Then \( A^{\text{op}} \) is an \( R \)-algebra with the same addition and \( R \)-action as \( A \) and with multiplication given by \( (xy)^{\text{op}} = y^{\text{op}}x^{\text{op}} \).

In particular, if \( M \) is an left \( A \)-module and a right \( B \)-module then \( M \) is a left \( A \otimes_F B^{\text{op}} \)-module via \( (a \otimes b^{\text{op}})m := amb \) with \( r(amb) = r(a \otimes b^{\text{op}})m = (ra \otimes b^{\text{op}})m = (a \otimes rb^{\text{op}})m = a(rm)b \), for all \( r \in R, a \in A, b \in B \), and \( m \in M \).

We also make the observation that if \( A \) is a simple \( R \)-algebra, then the \textit{center} of \( Z(A) := \{ x \in A \mid xy = yx \text{ for all } y \in A \} \) is a field as \( Ax = xA = AxA = A \), for all nonzero \( x \in Z(A) \).

Throughout this section \( F \) will denote a field.

**Definition 103.1.** Let \( F \) be a field and \( A \) a simple \( F \)-algebra. We say \( F \) is a \textit{central \( F \)-algebra} if \( F = Z(A) \). We say \( A \) is a \textit{finite dimensional \( F \)-algebra} if it is a finite dimensional \( F \)-vector space.

As mentioned above, the study of finite dimensional central simple \( F \)-algebras is bound to the study of finite dimensional central simple division \( F \)-algebras by Wedderburn’s Theorem 99.16. For example, we used this to investigate cyclic \( F \)-algebras, which are finite dimensional central simple \( F \)-algebras. We will also be interested when we extend the field \( F \). If \( K/F \) is a field extension and \( A \) an \( F \)-algebra, then the \( K \)-algebra \( K \otimes_F A \) will be denoted by \( A^K \). So if \( A \) is a finite dimensional \( F \)-algebra, \( \dim_K A^K = \dim_F A \). We begin with the following:

**Proposition 103.2.** Let \( F \) be a field, \( A \) and \( B \) two \( F \)-algebras. Then
(1) $Z(A \otimes_F B) = Z(A) \otimes_F Z(B)$.
(2) If $A$ is also central and simple and $B$ is simple, then $A \otimes_F B$ is a simple $F$-algebra.
(3) If both $A$ and $B$ are central and simple $F$-algebras, then so is $A \otimes_F B$.
(4) If $K/F$ is a field extension and $A$ is a finite dimensional central simple $F$-algebra, then so is $A^K$.

**Proof.** (1): Certainly $Z(A) \otimes_F Z(B) \subset Z(A \otimes_F B)$. Conversely, let $0 \neq z \in Z(A) \otimes_F Z(B)$. Write

$$z = \sum_{i=1}^{r} a_i \otimes b_i$$

with $a_i \in A$, $b_i \in B$ and $A = \{a_1, \ldots, a_r\}$, $B = \{b_1, \ldots, b_r\}$.

We may assume that $B$ is $F$-linearly independent. In particular, if $x \in A$, then

$$(x \otimes 1)z - z(x \otimes 1) = \sum_{i=1}^{r} (xa_i - a_ix) \otimes b_i.$$  

So $xa_i = a_ix$ for all $x \in B$, since $B$ is linearly independent. Hence, $z \in Z(A) \otimes B$.

We may now assume in (*) that we have chosen the $a_i \in Z(A)$ and $b_i \in B$ with $A$ an $F$-linearly independent set. The same argument then shows that all the $b_i$ lie in $Z(B)$.

(2): Let $0 \neq \mathfrak{A} \subset A \otimes_F B$ be an ideal. Choose

$$0 \neq x = \sum_{i=1}^{r} a_i \otimes b_i \in \mathfrak{A}, \quad \text{with} \quad a_i \in A, \ b_i \in B \text{ and } r \text{ minimal.}$$

It follows that $A = \{a_1, \ldots, a_r\}$ and $B = \{b_1, \ldots, b_r\}$ are $F$-linearly independent. Since $A$ is simple, $A = Aa_1A$, so there exists an equation $1 = \sum_{j=1}^{s} c_j a_1 d_j$ for some $c_j, d_j \in A$ and $s$. Set $a'_i = \sum_{j=1}^{s} c_j a_i d_j$ for $i = 2, \ldots, r$. Then

$$x_1 := \sum_{j=1}^{s} c_j x d_j = 1 \otimes b_1 + \sum_{i=2}^{r} a'_i \otimes b_i \text{ lies in } \mathfrak{A}.$$  

Since $B$ is linearly independent, $x_1$ is not zero. Applying the same argument to $Bb_1B$ produces $b'_i$ with $b'_1 = 1$ and

$$x_2 := 1 \otimes 1 + \sum_{i=2}^{r} a'_i \otimes b'_i \text{ lies in } \mathfrak{A}.$$
By the minimality of $r$, $A' = \{a'_1, \ldots, a'_r\}$ and $B' = \{b'_1, \ldots, b'_r\}$ are $F$-linearly independent. But if $a \in A$,
\[
ax_2 - x_2a = \sum_{i=2}^r (aa'_i - a'_ia) \otimes b'_i \text{ lies in } \mathfrak{A}.
\]
It follows that $aa'_i = a'_ia$ for all $a \in A$ by the minimality of $r$. In particular, $a \in Z(A)$. As $A$ is linearly independent, we must have $r = 1$.

Statements (3) and (4) follow from (1) and (2).

\begin{notation}
It is convenient to extend our notation $[K : F] = \dim_F K$ for field extensions $K/F$. In particular, if $A$ is a finite dimensional central $F$-algebra and $F \subset D \subset A$ a subalgebra with $D$ a division $F$-algebra, let $[A : D] := \dim_D A$, the $D$-dimension of the (left) $D$-vector space $A$. Of course, we have $[A : F] = [A : D][D : F]$.
\end{notation}

\begin{proposition}
Let $D$ be a finite dimensional central division $F$-algebra. Then $[D : F]$ is a square.
\end{proposition}

\begin{proof}
Let $\bar{F}$ be an algebraic closure of $F$. Then $D^{\bar{F}}$ is a finite dimensional central $\bar{F}$-algebra by Proposition 103.2. Since $\bar{F}$ is algebraically closed, $\bar{F}$ is the only finite dimensional division $\bar{F}$-algebra. Hence $D^{\bar{F}} \cong M_n(\bar{F})$, some $n$ by Wedderburn’s Theorem 99.16. Thus $[D : F] = [D^{\bar{F}} : \bar{F}] = n^2$.
\end{proof}

\begin{corollary}
Let $A$ be a finite dimensional central simple $F$-algebra. Then $[A : F]$ is a square.
\end{corollary}

\begin{proof}
By Wedderburn’s Theorem 99.16, $A \cong M_n(D)$ for some finite dimensional division $F$-algebra $D$. So $[A : D]$ is a square. As $Z(M_n(D)) = Z(D)$, we must have $Z(D) = F$, i.e., $D$ is central. Since $[A : D]$ is a square by the previous result and $[A : F] = [A : D][D : F]$, the result follows.
\end{proof}

The corollary allows us to define an important invariant of finite dimensional central simple $F$-algebras.

\begin{definition}
Let $A$ be a finite dimensional central simple $F$-algebra, so $[A : F]$ is a square. The \textit{degree} of $A$ is defined by $\deg A := \sqrt{[A : F]}$.
\end{definition}

\begin{definition}
Let $A$ be an $R$-algebra and $B \subset A$ a subalgebra. We let $Z_A(B) := \{x \in A \mid xb = bx \text{ for all } b \in B\}$, the \textit{centralizer} of $B$ in $A$. This is an $R$-algebra.
\end{definition}

We can now prove the following basic theorem:
**Theorem 103.8.** (Double Centralizer Theorem) Let $A$ be a finite dimensional central simple $F$-algebra and $B$ a simple subalgebra of $A$. Then the following are true:

1. $Z_A(B)$ is simple.
2. $B = Z_A(Z_A(B))$.
3. $[A : F] = [B : F][Z_A(B) : F]$.
4. If $B$ is a finite dimensional central simple $F$-algebra, then $Z_A(B)$ is a finite dimensional central simple $F$-algebra and $A = B \otimes_F Z_A(B)$.

**Proof.** Let $C = Z_A(B)$ and $T = B \otimes_F A^{\text{op}}$. By Proposition 103.2, $T$ is a simple algebra, as $A^{\text{op}}$ is clearly one. Let $\mathfrak{A}$ be a minimal left ideal of $T$ and $D = \text{End}_T(\mathfrak{A})$, a division ring by Schur’s Lemma 99.15. The $F$-algebra $A$ becomes a $T$-module by $(b \otimes a^{\text{op}})a = bxa$ for all $a, x \in A$, $b \in B$. We determine $\text{End}_T(A)$. If $c \in C = Z_A(B)$, then $\lambda_c : A \to A$ given by $a \mapsto ca$ is $F$-linear as

$$
\lambda_c((b \otimes a^{\text{op}})(x)) = cbxa = bcxa = (b \otimes a^{\text{op}})(cx) = (b \otimes a^{\text{op}})\lambda_c(x),
$$

so lies in $\text{End}_T(A)$. Next suppose that $f \in \text{End}_T(A)$. Then

$$(*) \quad f(a) = f((1 \otimes a^{\text{op}})(1) = (1 \otimes a^{\text{op}})f(1) = \lambda_{f(1)}(a).$$

It follows that $\text{End}_T(A) = \{\lambda_{f(1)} \mid f \in \text{End}_T(A)\}$. Let $f \in \text{End}_T(A)$. We show $f(1) \in C$. Let $b \in B$. As $B \subset A$, we have

$$bf(1) = (b \otimes 1^{\text{op}})f(1) = f((b \otimes 1^{\text{op}})1) = f(b) = f(1)b.$$

Therefore, $f(1) \in Z_A(B) = C$, and $\text{End}_T(A) = \{\lambda_c \mid c \in C\}$.

Now $A$ is a finite dimensional $F$-algebra, so $A$ is a finitely generated $T$-module. As $T$ is a simple algebra, $A \cong \mathfrak{A}^n$ for some $n$ as a $T$-module, hence

$$C \cong \text{End}_T(A) \cong \text{End}_T(\mathfrak{A}^n) \cong \mathcal{M}_n(D)$$

is simple. This proves (1).

We also have that $T \cong \mathfrak{A}^r$ some $r$, so $T \cong \mathcal{M}_r(D)$ and $\mathfrak{A} \subset T$ corresponds to a column space of $\mathcal{M}_r(D)$, i.e., $[\mathfrak{A} : D] = r$. Thus

$$[A : F] = [A : D][D : F] = nr[D : F]$$

$$[A : F][B : F] = [T : F] = [\mathcal{M}_r(D) : F] = r^2[D : F]$$

$$[C : F] = n^2[D : F].$$

Therefore,


Statement (3) now follows.
Finally, we show (4). Clearly, \( B \subset Z_A(Z_A(B)) = Z_A(C) \). By (1), \( C \) is simple, so applying (3) to \( C \subset A \) yields

\[
[A : F] = [C : F][Z_A(C) : F] = [C : F][B : F].
\]

So \([Z_A(C) : F] = [B : F]\). It follows that \( B = Z_A(C) \). If \( B \) is a finite dimensional central simple \( F \)-algebra, then \( B \otimes_F C \) is simple by Proposition 103.2 and \([B \otimes_F C : F] = [B : F][C : F] = [A : F]\), so the map \( B \otimes_F A \rightarrow A \) induced by \( b \otimes c \mapsto bc \) is an \( F \)-algebra isomorphism with

\[
F = Z(A) \cong Z(B \otimes_F C) = Z(B) \otimes_F Z(C) = F \otimes_R Z(C) \cong Z(C).
\]

Statement (4) now follows. \( \square \)

We leave as an easy exercise:

**Remarks 103.9.** Let \( A \) be an \( R \)-algebra. Then we have \( R \)-algebra isomorphisms:

1. \( \mathcal{M}_n(R) \otimes_R A \cong \mathcal{M}_n(A) \cong A \otimes \mathcal{M}_n(R) \).
2. \( (\mathcal{M}_n(A))^\text{op} \cong \mathcal{M}_n(A^\text{op}) \).
3. \( \mathcal{M}_n(R) \otimes_R \mathcal{M}_m(R) \cong \mathcal{M}_{nm}(R) \).

**Corollary 103.10.** Let \( A \) be a finite dimensional central simple \( F \)-algebra and \( B \subset A \) a simple subalgebra. Then

1. \( B \otimes_F A^\text{op} \cong \mathcal{M}_s(Z_A(B)) \) with \( s = [B : F] \).
2. \( A \otimes_F A^\text{op} \cong \mathcal{M}_m(F) \) with \( m = [A : F] \).

**Proof.** In the notation set up in the proof of the Double Centralizer Theorem, we have

\[
T = B \otimes_R A^\text{op} \cong \mathcal{M}_r(D) \quad \text{and} \quad Z_A(B) \cong \mathcal{M}_n(D).
\]

We also know that \([A : F] = [B : F][Z_A(B) : F]\). so if \( s = [B : F] \), we have

\[
r^2[D : F] = [T : F] = [B : F][A : F]
\]

\[
= [B : F]^2[Z_A(B) : F] = s^2n^2[D : F].
\]

If follows that \( r = sn \) and

\[
B \otimes_R A^\text{op} \cong \mathcal{M}_r(D) \cong \mathcal{M}_{sn}(D) \cong \mathcal{M}_s(F) \otimes_F \mathcal{M}_n(F) \otimes_F D
\]

\[
\cong \mathcal{M}_s(F) \otimes_F \mathcal{M}_n(D) \cong \mathcal{M}_s(F) \otimes_F Z_A(B) \cong \mathcal{M}_s(Z_A(B)).
\]

This establishes (1) and (2) follows from (1). \( \square \)

We next turn to subfields in a division \( F \)-algebra. We begin with some definitions.
Definition 103.11. Let $A$ be an $F$-algebra and $K/F$ a field extension. We say

1. We say that $A$ is split over $F$ if $A \cong \mathbb{M}_n(F)$ some integer $n$.
2. We say that $A$ is $K$-split or splits over $K$ if $A^K$ is split. If this is the case, we call $K$ a splitting field for $A$.
3. We say $K$ is a maximal subfield of $A$ if $K \subset A$ and is a maximal such.
4. If $A \cong \mathbb{M}_n(D)$, with $D$ a finite dimensional division $F$-algebra, we let $\text{ind} A := \deg D = \sqrt{|D : K|}$ called the (Schur) index of $A$.

We now show that there exist subfields of a finite dimensional division $F$-algebra that split $A$. We first note the following:

Remarks 103.12. 1. Let $A$ be a finite dimensional central simple $F$-algebra. Then $\text{ind} A | \deg A$.
2. If $A \cong \mathbb{M}_n(D)$ with $D$ a finite dimensional division $F$-algebra and $K/F$ a field extension, then $K$ splits $A$ if and only if it splits $D$.

Corollary 103.13. Let $D$ be a finite dimensional division $F$-algebra and $K \subset D$ be a maximal subfield. Then $[K : F] = [D : K] = \deg D$.

Proof. Since $K$ is simple, $[D : F] = [K : F][Z_D(K) : F]$ by the Double Centralizer Theorem. We observe that the maximal subfield $K \subset D$ must satisfy $K = Z_D(K)$. For if not, there exists an $a \in Z_D(K) \setminus K$ with $K < K(a)$ a subfield of $D$, a contradiction. Therefore, $[K : F][K : F] = [D : F] = [D : K][K : F]$, by Proposition 103.4. The result follows.

Corollary 103.14. Let $A$ be a finite dimensional central simple $F$-algebra. If $A \cong \mathbb{M}_n(D)$, $D$ a division $F$-algebra, $K$ a maximal subfield of $D$, then ind $A = [K : F]$.

If $A$ is a finite dimensional central simple $F$-algebra, it is possible for a maximal subfield $K \subset A$ to satisfy $\deg A > [K : F]$.

Corollary 103.15. Every maximal subfield $K$ of a finite dimensional division $F$-algebra $D$ splits $D$.

Proof. We know that $Z_D(K) = K$, so for a unique integer $s$, $K \otimes_F D^{op} \cong \mathbb{M}_s(K)$ some $s$ by Corollary ???. Hence

$$K \otimes_F D = (K \otimes_F D^{op})^{op} \cong (\mathbb{M}_s(K))^{op} \cong \mathbb{M}_s(K)$$

as needed.

In order to use Galois Theory in the study of finite dimensional central simple $F$-algebras, we need the following crucial proposition.
Proposition 103.16. Let $D$ be a finite dimensional division $F$-algebra. Then $D$ contains a maximal subfield that is finite and separable over $F$. In particular, $D$ contains a separable splitting field.

Proof. We may assume that $\text{char } F = p > 0$. We induct on $n = \text{ind } D = \deg D = \sqrt{|D : F|}$. We may assume that $n > 1$.

Claim. There exists $F < E \subset D$ with $E/F$ a separable field extension.

Suppose that this is false. We first show that every element $x \in D$ is purely inseparable over $F$, i.e., $x^{p^e}$ lies in $F$ some integer $e$. Indeed if $x \in D$, then the minimal polynomial $m_F(x) = f(t)$ in $F[t]$ for some separable polynomial $f \in F[t]$. Our assumption implies that $f$ splits over $F$, hence $x^{p^e}$ lies in $F$, and every element in $D$ is purely inseparable over $F$. In particular, by Proposition 103.4, this means that for all $x \in D \setminus F$, we have $n^2 = [D : F] = [D : F(x)] [F(x) : F]$. Hence $p | n$.

Let $K$ be a maximal subfield of $D$. By the last corollary, $K$ splits $D$, so there exists a $K$-algebra isomorphism $\varphi : K \otimes_F D \to \mathfrak{m}_n(K)$. Let $x \in D \setminus F$. As $x$ is purely inseparable over $F$, $a = x^{p^e}$ lies in $F$ for some $e$. In particular,

$$\varphi(1 \otimes x)^{p^e} = \varphi(1 \otimes x^{p^e}) = \varphi(1 \otimes a) = aI,$$

(where $I$ is the identity matrix). Let $\tilde{K}$ be an algebraic closure of $K$. Then $(\varphi(1 \otimes x) - a^{1/p^e} I)^{p^e} = 0$. In particular, all the eigenvalues of $\varphi(1 \otimes x)$ must be equal. Since $p | n$, the trace $\text{tr } (\varphi(1 \otimes x)) = 0$ in $K$ (since in $\tilde{K}$). But the set $\{ \alpha \in M_n(K) | \text{tr}(\alpha) = 0 \}$ cannot span $M_n(K)$ as the trace is $K$-linear. This is a contradiction and establishes the claim.

Let $F < E \subset D$ with $E/F$ separable and $D' = Z_D(E)$ the subdi- vision $F$-algebra of $D$. We next show that $E = Z(D')$. By the Double Centralizer Theorem 103.8, $E = Z_D(D')$, so $Z(D') \subset E$. Conversely, if $x \in E$, then $D' = Z_D(E)$ means that $xd' = d'x$ for all $d' \in D'$. Thus $E = Z_D(D')$. It follows that $D$ is a finite dimensional division $E$-algebra. By the Double Centralizer Theorem 103.8, we have

$$[D : F] = [D' : F][Z_D(D') : F] = [D' : F][E : F] = [D' : E][E : F]^2 > [D' : E]$$

by the claim. Hence by induction, there exists a maximal subfield $L$ of $D'$ with $L/E$ separable. In particular, $L/F$ is also separable. To finish it suffices to show that $L$ is a maximal subfield of $D$. If this is not the case, then $L < Z_D(L)$. But $Z_D(L) \subset Z_D(E) = D'$. This contradicts the maximality of $L$. □
Corollary 103.17. Let $A$ be a finite dimensional central simple $F$-algebra. Then $A$ has a separable splitting field finite over $F$. In particular, $A$ also has a Galois splitting field finite over $F$.

**Proof.** Let $A \cong M_n(D)$ with $D$ a finite dimensional division $F$-algebra. Then $D$ contains a maximimal subfield $K$ that is separable over $F$. Thus $K$ splits $D$, hence $A$. Let $L/K$ be the normal closure of $K/F$. Then $L/F$ is Galois and $L$ splits $A$. $\square$

Perhaps the most important basic theorem is the following:

**Theorem 103.18.** (Skolem-Noether Theorem) Let $A$ be a simple, left Artinian central $F$-algebra and $B \subset A$ a finite dimensional simple $F$-subalgebra. If $f, g : B \to A$ are $F$-algebra maps, then there exists a unit $u \in A^\times$ satisfying

$$f(b) = u^{-1}g(b)u \text{ for all } b \in B.$$

**Proof.** By Proposition 103.2, $C = B \otimes_F A^{\text{op}}$ is a simple $F$-algebra, and by Wedderburn’s Theorem 99.16, there exists a division $F$-algebra $D$ such that $A \cong M_n(D)$ for some $n$. Then

$$C \cong B \otimes_F (M_n(D))^{\text{op}} \cong B \otimes_F M_n(D^{\text{op}}) \cong M_n(B \otimes_F D^{\text{op}})$$

is left Artinian, since it is a finite dimensional $D^{\text{op}}$-algebra with $D^{\text{op}}$ a division ring. Let $M$ be an irreducible $C$-module, so unique up to isomorphism with $C$-action induced by $(b \otimes a^{\text{op}})(m) = bma$ for all $b \in B$, all $a \in A$, and all $m \in M$. Define $C$-modules $A_1$ and $A_2$ as follows: $A_1 = A$ with $C$-action induced by

$$(b \otimes a^{\text{op}})(x) = f(b)xa \text{ for all } b \in B \text{ and all } a, x \in A_1$$

and $A_2 = A$ with $C$-action induced by

$$(b \otimes a^{\text{op}})(x) = g(b)xa \text{ for all } b \in B \text{ and all } a, x \in A_2.$$ 

Then there exist unique integers $n_i$ such that $A_i \cong M^{n_i}$ as $C$-modules, since $A_i$ is a finitely generated $C$-module for $i = 1, 2$. Since $A$ is a left $D$-vector space, $M$ is a right $D$-vector space. Therefore,

$$n_1 = \frac{[A : D]}{[M : D]} = n_2$$

as $[A_i] = [A : D]$ for $i = 1, 2$. It follows that $A_1 \cong A_2$ as $C$-modules.

Let $\varphi : A_1 \to A_2$ be a $C$-isomorphism. Then we have $\varphi((b \otimes a^{\text{op}}))(x) = (b \otimes a^{\text{op}})\varphi(x)$ for all $b \in B, a, x \in A$, i.e.,

$$\varphi((f(b)xa) = g(b)\varphi(x)a \text{ for all } a, x \in A, b \in B.$$

Let $u = \varphi(1)$. Setting $b = 1 = x$ in (*) shows that $f(1) = 1 = g(1)$.

Consequently, $\varphi(x) = \varphi(1)x = ux$ for all $x \in A$. As $\varphi$ is surjective,
there exists a $v \in A$ satisfying $uv = 1$. Since $A \cong \mathbb{M}_n(D)$, we have $u \in A^\times$. Setting $x = 1 = a$ in (*), we see that $uf(b) = \varphi(1 \cdot f(b)) = \varphi(f(b) \cdot 1) = g(b)u$ for all $b \in B$ as $f(b) \in A$ and $\varphi$ is $C$-linear. The result follows.

**Corollary 103.19.** Let $A$ be a finite dimensional central simple $F$-algebra. Then every $F$-algebra endomorphism of $A$ is an inner $A$-automorphism.

**Proof.** Any $A$-algebra homomorphism takes $1_A$ to $1_A$, so is not trivial.

**Corollary 103.20.** Let $A$ be a finite dimensional central simple $F$-algebra. They any two isomorphic simple subalgebras of $A$ are conjugate and, in particular, have conjugate centralizers.

As an application of Skolem-Noether, we give an alternative proof of Wedderburn’s Theorem about the commutativity of finite division rings.

**Lemma 103.21.** Let $D$ be a finite dimensional division $F$-algebra with $K$ a maximal subfield of $D$. Then all maximal subfields of $D$ are isomorphic to $K$ as $F$-algebras if and only if $K/F$ is separable and $D = \bigcup_{x \in D^\times} x^{-1}Kx$.

**Proof.** ($\Rightarrow$): By Proposition 103.16, $D$ contains a separable maximal subfield. As all maximal subfields are isomorphic by assumption, all maximal subfields are separable over $F$. Since every $x \in D$ lies in some maximal subfield $E$ and $E \cong K$ as $F$-algebras, we have $E = x^{-1}Kx$ for some $x \in D^\times$ by the Skolem-Noether Theorem. The result follows.

($\Leftarrow$): Let $y \in D$. By hypothesis, there exists an element $x \in D^\times$ such that $y \in x^{-1}Kx$. Thus $y$ is separable over $F$ for all $y \in D^\times$, hence all maximal subfields in $D$ are separable. If $L$ is a maximal subfield of $D$, then $L = F(y)$ for some $y \in D^\times$ by the Primitive Element Theorem 56.9 and $y \in x^{-1}Kx$, for some $x \in D^\times$. Thus $L = F(y) = x^{-1}Kx$ by maximality.

**Theorem 103.22.** (Wedderburn) Every finite division ring $D$ is a field.

**Proof.** Let $F = Z(D)$, so $D$ is a finite dimensional central division $F$-algebra, say of degree $n$. Then all maximal subfields of $D$ have $|F|^n$ elements, and there is precisely one up to isomorphism, say $K$. By the lemma, $D^\times = \bigcup_{x \in D^\times} x^{-1}K^\times x$. If $D^\times \neq E^\times$, then the number of conjugates $x^{-1}K^\times x$ of $K^\times$, $x \in D^\times$, is at most $[D^\times : F^\times]$. Since $1 \in
for all \( x \in D^\times \), we have \(|\{x^{-1}hx \mid x \in D^\times, h \in K^\times\}| < D^\times\), a contradiction. So \( D = K \), hence \( D = F \). \( \square \)

We saw in Section 98 that generalized quaternion algebras have norms attached to them and are useful in deciding whether they were division rings or matrix rings. Using the Skolem-Noether Theorem, we can define a norm called the reduced norm on any finite dimensional central simple \( F \)-algebra. We now construct it.

**Construction 103.23.** Let \( A \) be a finite dimensional central simple \( F \)-algebra. We know by Corollary 103.17 that \( A \) has a separable splitting field \( K \), finite over \( F \). Let \( \deg A = n \) and \( \varphi : A^K \to \mathbb{M}_n(K) \) be an \( K \)-algebra isomorphism. If \( \alpha \in \mathbb{M}_n(F) \), let \( f_\alpha := \det(tI - \alpha) \) denote the characteristic polynomial of \( \alpha \) in \( K[t] \). If \( a \in A \), define the reduced characteristic polynomial of \( a \) by \( f_{\varphi(a)} \in K[t] \) over \( K \). This is independent of the isomorphism \( \varphi \) by the Skolem-Noether Theorem 103.18. If \( f_{\varphi(a)} = t^n + c_{n-1}t^{n-1} + \cdots + c_0 \in K[t] \), define the reduced norm of \( a \) over \( K \) by \( \text{Nrd}(a) := (-1)^nc_0 \) and the reduced trace of \( a \) over \( K \) by \( \text{Trd}(a) := -c_0 \) over \( K \).

We want to show that the reduced norm of a finite dimensional central simple \( F \)-algebra is independent of the separable splitting field \( K \) and takes values in \( F \). We begin by looking at the reduced characteristic polynomial.

**Lemma 103.24.** Let \( A \) be a finite dimensional central simple \( F \)-algebra of degree \( n \). Then the reduced characteristic polynomial of an element in \( A \) is independent of the finite separable splitting field of \( A \) over \( F \). In particular, we can evaluate the reduced characteristic polynomial of any element using a finite Galois extension of \( F \) that splits \( A \).

**Proof.** The algebra \( A \) has a separable splitting field \( K \), finite over \( F \). Let \( \tilde{F} \) be an algebraic closure of \( F \) containing \( K \). If \( L \) is another separable splitting field of \( A \) over \( F \) in \( \tilde{F} \), then we have \( E = L(K) = K(L) \) is also a finite separable splitting field of \( A \). It follows that the reduced characteristic polynomials of \( \alpha \in A \) over \( \tilde{K} \) and over \( L \) are the same as they are the same over \( \tilde{F} \), since \( \mathbb{M}_n(E) \subset \mathbb{M}_n(\tilde{F}) \) for all intermediate fields \( \tilde{F}/E/F \). Therefore, the reduced characteristic polynomial of \( a \) is independent of separable splitting field. In particular, we may assume that \( L/F \) is finite Galois with \( L/K/F \). \( \square \)

We now look at the action of a finite Galois group of a Galois splitting field on a finite dimensional central simple \( F \)-algebra.

**Proposition 103.25.** Let \( A \) be a finite dimensional central simple \( F \)-algebra of degree \( n \). Then \( f_{\varphi(a)} \in F[t] \) for all \( a \in A \).
Proof. Let $K/F$ be a finite separable field extension with $A$ split over $K$. Let $L/K/F$ be a finite field extension with $L/F$ is Galois. For all $\sigma \in G(L/F)$, the map $\sigma \otimes 1_A : A^L \to A^L$ induced by $x \otimes a \mapsto \sigma(x) \otimes a$ for $x \in L$, and $a \in A$ is a $K$-algebra isomorphism. Let the map $\mathbb{M}_n(L) \to \mathbb{M}_n(L)$ be defined by taking the matrix $(x_{ij})$ to $(\sigma(x_{ij}))$. Denote this map by $\sigma$ also. We have a commutative diagram

$$
\begin{array}{ccc}
A^L & \xrightarrow{\sigma \otimes 1_A} & A^L \\
\varphi \downarrow & & \downarrow \psi \\
\mathbb{M}_n(L) & \xrightarrow{\sigma} & \mathbb{M}_n(L)
\end{array}
$$

with $\psi = \sigma \varphi (\sigma \otimes 1_A)^{-1}$, a ring isomorphism. Since $\psi(x \otimes 1) = \sigma \varphi (\sigma^{-1}(x) \otimes 1) = \begin{pmatrix} x & \cdots \\ \vdots \\ x \end{pmatrix}$ for all $x \in L$, it follows that $\psi$ is an $L$-isomorphism. If $a \in A$, we have $\psi(1 \otimes a) = \sigma \varphi ((\sigma \otimes 1_A)^{-1})(1 \otimes a) = \sigma(\varphi(1 \otimes a))$.

By the Skolem-Noether Theorem 103.18, $\varphi(1 \otimes a)$ and $\psi(1 \otimes a)$ have the same characteristic polynomials, so its coefficients lie in $F$. The result follows.

Corollary 103.26. Let $A$ be a finite dimensional central simple $F$-algebra of degree $n$. Then for all $a, b \in A$ and all $x \in F$, we have

1. The reduced norm of $A$ satisfies $\text{Nrd} : A \to F$. Moreover, $\text{Nrd}(ab) = \text{Nrd}(a) \text{Nrd}(b)$ and $\text{Nrd}(xa) = x^n \text{Nrd}(a)$.
2. An element $a \in A$ is a unit if and only if $\text{Nrd}(a)$ is nonzero. In particular, $\text{Nrd} : A^\times \to F^\times$ is a group homomorphism.
3. The reduced trace of $A$ satisfies $\text{Trd} : A \to F$ is $F$-linear. Moreover, $\text{Trd}(ab) = \text{Trd}(ba)$.

Proof. By the proposition, all the coefficients of $f_a, a \in A$ lie in $F$. In particular, $\text{Nrd}(a)$ and $\text{Trd}(a)$ lie in $F$ for all $a \in A$. The other statements follow as the same properties hold for the corresponding properties for the determinant and trace of matrices.

Remarks 103.27. Let $A$ be a finite dimensional central simple $F$-algebra of degree $n$ and $K/F$ a finite separable field extension with $A$ split over $K$. Let $\mathcal{B} = \{a_1, \ldots, a_n\}$ be a basis for $A$ and $\varphi : A^K \to \mathbb{M}_n(K)$ an $F$-algebra isomorphism. Then $\text{Nrd}(\sum_{i=1}^n x_i a_i) = \det (\sum_{i=1}^n x_i \varphi(a_i))$ defines a polynomial $\text{Nrd}_A \in F[t_1, \ldots, n^2]$. As we
are taking a determinant, this is a homogeneous polynomial of degree $n$.

We arrive at the result that we sought:

**Theorem 103.28.** (Tsen) Let $F$ be a field of transcendence degree one over an algebraically closed field. Then every finite dimensional central simple $F$-algebra is isomorphic to a matrix ring over $F$.

**Proof.** If $A$ is a finite dimensional central simple $F$-algebra of degree $n$, then the norm form $\text{Nrd}_A$ is a homogeneous polynomial of degree $n$. As $F$ is a $C_1$-field by Corollary 98.17, it has a nontrivial zero if $n > 1$, hence cannot be a matrix ring over a noncommutative division $F$-algebra. □

**Exercises 103.29.**

$F$ is a field in all the exercises.

1. Let $A$ be a finite dimensional $F$-algebra and $K/F$ a finite field extension. Prove that $A$ is a finite dimensional central simple $F$-algebra if and only if $A^K$ is a finite dimensional central simple $F$-algebra.

2. Let $A$ be a finite dimensional $F$-algebra. Show that $A$ is a finite dimensional central simple $F$-algebra if and only if there exists a finite extension $K/F$ such that $A$ splits over $K$.

3. Let $\mathfrak{A}$ be a nonzero left ideal in a finite dimensional simple $F$-algebra $A$. Set $D = \text{End}_A(\mathfrak{A})$. Show the map $\lambda_F : A \to \text{End}_D(\mathfrak{A})$ defined by $a \mapsto \lambda_a : x \mapsto ax$ is an ring isomorphism.

4. The $F$-subalgebras in $M_n(F)$ that are isomorphic to $F^n$ are conjugate to the subalgebra of diagonal matrices.

5. Show a finite dimensional central simple $F$-algebra $A$ of degree $n$ is split if and only if it contains a $F$-subalgebra isomorphic to the direct product $F^n = F \times \cdots \times F$.

6. Let $D$ be a central division $F$-algebra. If $a, b \in D$ are algebraic over $F$ with the same minimal polynomial over $F$, show that $b = xax^{-1}$ for some $x \in D^\times$.

7. Let $D$ be a central division $F$-algebra with $\text{char } F = p > 0$. Let $a \in D \setminus F$ satisfy $a^{p^n} = a$ for some positive integer $n$. Show that there exists an element $x \in D^\times$ satisfying $xax^{-1} = a^i \neq a$ for some $i$.

8. Let $A$ be an $F$-algebra. An $F$-derivation $\delta : A \to A$ is an $F$-linear map satisfying $\delta(ab) = a\delta(b) + \delta(a)b$. Suppose that $A$ is a central simple $F$-algebra and $\delta : A \to A$ is an $F$-derivation. Show the following:
(a) The subalgebras 

\[ B = \left\{ \begin{pmatrix} a & \delta(a) \\ 0 & a \end{pmatrix} \mid a \in A \right\} \quad \text{and} \quad B = \left\{ \begin{pmatrix} a & 0 \\ 0 & a \end{pmatrix} \mid a \in A \right\} \]

of \( \mathbb{M}_2(A) \) are isomorphic \( F \)-algebras.

(b) \( \delta : A \to A \) is an \( F \)-inner derivation, i.e., there exists an element \( c \in A \) such that \( \delta(x) = xc - cx \) for every \( x \in A \).

9. Let \( D \) be a finite dimensional central division \( F \)-algebra with \( F \) a \( C_2 \)-field. Show \( \text{Nrd} : D \to F \) is surjective.

104. The Brauer Group

To further study the theory of finite dimensional central simple \( F \)-algebras, it is useful to give a structure arising from the set of isomorphism classes of such algebras. We shall show that the tensor product of finite dimensional central simple \( F \)-algebras is also a finite dimensional central simple \( F \)-algebra. Since any finite dimensional central simple \( F \)-algebra is a matrix ring over a central division \( F \)-algebra, unique up to isomorphism, the question arises about the underlying division ring in the tensor product of two such central simple algebras. We construct a group called the Brauer group of \( F \) group whose group structure is induced by the tensor of such algebras. This group is a torsion group and is an important invariant in the study of fields. Historically, it was crucial in the interpretation of parts of number theory (class field theory).

Construction 104.1. Let \( A_i \) be a finite dimensional central simple \( F \)-algebra with \( A_i \cong \mathbb{M}_{n_i}(D_i) \), \( D_i \) a central division \( F \)-algebra for \( i = 1, 2 \). We say that \( A_1 \) is (Brauer) equivalent to \( A_2 \) if \( D_1 \cong D_2 \) as \( F \)-algebras. Let \([A_i]\) denote the (Brauer) equivalence class of \( A_1 \). Since \( A_1 \otimes_F A_2 \) is a finite dimensional central simple \( F \)-algebra and \( A_1 \otimes_F A_1^{op} \cong \mathbb{M}_{n_1^2}(F) \), we have

\[ \text{Br}(F) := \{ [A] \mid A \text{ is a finite dimensional central simple } F\text{-algebra} \} \]

is an abelian group via \([F] = 1\), \([A_1][A_2] := [A_1 \otimes_F A_2]\), and \([A_1]^{-1} = [A_1^{op}]\), called the Brauer Group of \( F \). Let \( K/F \) be a field extension. Then

\[ (A_1 \otimes_F A_2)^K = K \otimes_F A_1 \otimes_F A_2 \cong (K \otimes_F A_1) \otimes_K (K \otimes_F A_2) \cong A_1^K \otimes_F A_2^K \]

and

\[ (A_1)^{op} = (K \otimes_F A_1)^{op} \cong K \otimes_F A_1^{op} = (A_1^{op})^K. \]
Therefore, we have a group homomorphism $\varphi : \text{Br}(F) \to \text{Br}(K)$ given by $[A] \mapsto [A^K]$. Let

$$\text{Br}(K/F) := \ker \varphi = \{[A] | A \text{ is a finite dimensional central simple } F\text{-algebra with } A^K \text{ split}\}.$$ 

Let $\bar{F}$ be a fixed algebraic closure of $F$ and $F_{\text{sep}}$ denote the separable closure of $F$ in $\bar{F}$, i.e., the maximal separable extension of $F$ in $\bar{F}$. The Galois group of $G(\bar{F}/F) = G(F_{\text{sep}}/F)$ called the absolute Galois group of $F$. Since every finite dimensional central simple $F$-algebra has a finite separable splitting field, we saw that it has a finite Galois splitting field. It follows that

$$\text{Br}(F) = \bigcup_{K/F \text{ finite separable}} \text{Br}(K/F) = \bigcup_{L/F \text{ finite Galois}} \text{Br}(L/F),$$

and $\text{Br}(F) = \text{Br}(F_{\text{sep}}/F)$.

**Proposition 104.2.** Let $D$ be a finite dimensional central division $F$-algebra with $\deg D = q$ and $K/F$ a finite field extension. Then

1. If $D$ splits over $K$, then $q \mid [K : F]$.
2. There exists a smallest positive integer $r$ such that $K$ can be embedded in $M_r(F)$.
3. Let $r$ be as in (2). Then $D$ splits over $K$ if and only if $K$ is a maximal subfield of $M_r(D)$.
4. Let $r$ be as in (2) and $E = Z_{M_r(D)}(K)$. Then $E$ is a division $F$-algebra. Moreover, $K$ is a maximal subfield of $M_r(D)$ if and only if $E = K$.

**Proof.** Let $A = D^K$, a finite dimensional central simple $F$-algebra, and $M$ an irreducible $A$-module. Then $M$ is a $D$-vector space by $dm = (1 \otimes d)m$ for all $d \in D, m \in M$. Let $r = [M : D]$. We have

$$(k \otimes 1)(1 \otimes d) = k \otimes d = (1 \otimes d)(k \otimes 1) \text{ for all } k \in K, d \in D.$$ 

It follows that $\lambda_k : M \to M$ determined by $m \mapsto (k \otimes 1)m$ is $D$-linear for all $k \in K$, hence $K = K \otimes_F F \subset \text{End}_D(M) = M_r(D)$, where we view $M$ as a right $\text{End}_D(M)$-module. [We do this so that we do not get $M_r(D^{\text{op}})$.]

**Claim 1.** The above $r$ satisfies (2):

Suppose that $K \subset M_s(D)$. Choose $N$ a $D$-vector space of dimension $s$ such that $K \subset \text{End}_D(N) \cong M_s(D)$, where we write linear operators on the right, matrices as usual. Let $N$ be an $A$-module via $(k \otimes d)v := dv\lambda_k$
for all \( k \in K, \ d \in D, \) and \( v \in N \). Since \( A \) is a finite dimensional central simple \( F \)-algebra, \( M \) is irreducible, so \( N \cong M^n \) for some \( n \), hence
\[
\text{End}_D(N) \cong \text{End}_D(M^n) \cong M_{rn}(D), \quad \text{i.e.,} \quad r \mid s.
\]
This proves Claim 1.

Let \( B = \text{End}_B(M) \) and \( E = Z_B(K) \), so \( [B : D] = r^2 \).

**Claim 2.** \( E = \text{End}_A(M) \):

Let \( k \in D, \ m \in M, \) and \( f \in \text{End}_A(M) \). Then
\[
mf\lambda_k = (d \otimes 1)(mf) = (k \otimes 1)m f = m\lambda_k f,
\]
so \( f \in E \). Conversely, if \( f \in E \), then for all \( k \in K, \ m \in M, \) and \( d \in D, \)
\[
((k \otimes d)m) f = ((k \otimes 1)(1 \otimes d)f = ((1 \otimes d)m\lambda_k)f \\
= (1 \otimes d)(mf)\lambda_k = (k \otimes 1)(1 \otimes d)(mf) = (k \otimes d)(mf),
\]
so \( f \in \text{End}_A(M) \) and Claim (2) is established.

We next show (4) holds. By Schur’s Lemma, \( E \) is a division algebra. Suppose that \( K = E \). If \( K \subset L \subset B \), with \( L \) a field, then \( L \subset Z_B(K) = E = K \). Therefore, \( K \) is a maximal subfield of \( B \). Conversely, if \( K < E \), then there exits an element \( x \in E \setminus K \), so \( K < K(x) \subset B \). But \( K(x) \) is a field, so \( K \) is not a maximal subfield of \( B \), a contradiction. This establishes (4).

We have \( q = \deg D = \sqrt{[D : F]} \) and \( [B : F] = [B : D][D : F] \). So by the Double Centralizer Theorem 103.8
\[
(*) \quad [K : F][E : F] = [B : F] = r^2 q^2.
\]

We now show (1). suppose that \( K \) splits \( D \), hence \( A = D^K = \mathbb{M}_q(K) \) and \( A \cong M^q \) as an \( A \)-module (as \( M \) corresponds to a column space of \( \mathbb{M}_q(K) \). Therefore,
\[
= [A : K][K : F] = [D^K : K][K : F] = [D : F][K : F].
\]

Therefore,
\[
(\dagger) \quad qr = [A : D] = [K : F].
\]
which is (1).

Lastly, we prove (3).

\((\Rightarrow)\): We have \( K \subset E \) and by (\dagger), \( qr = [K : F] \). Therefore, by (\(\ast\)), we have \( [E : F] = qr \). It follows that \( E \cong K \).
Suppose that $K$ is a maximal subfield of $B$. We must show that $D$ splits over $K$. We know that $K = E = Z_B(K)$ by (4). By Corollary 103.10 of the Double Centralizer Theorem 103.8, we know that 

$$K \otimes K B^\text{op} \cong M_s(Z_B(K)) = M_s(E) = M_s(K)$$

with $s = [K : F]$. Hence $K \otimes B \cong M_s(K)$ also. Consequently, $B$ splits over $K$, hence so does $D$. \hfill \box

**Remark 104.3.** Let $D$ be a finite dimensional division $F$-algebra and $K/F$ a finite field extension. Suppose that $D$ splits over $K$. Then the proposition above shows that there exists a finite dimensional central simple $R$-algebra $B$ such that $[B] = [D]$ in the Brauer group $\text{Br}(F)$ satisfying

1. $K \subset B$ is a maximal subfield.
3. $K = Z_B(K)$.

We call such a $K$ a *self centralizing maximal subfield* of $B$. [If $B$ is not a division $F$-algebra, a maximal subfield $L$ of $B$ may not satisfy $Z_B(L) = L$.]

**Corollary 104.4.** Let $K/F$ be a field extension and $A$ a finite dimensional central simple $F$-algebra satisfying $[K : F]^2 = [A : F]$. Suppose that $A$ splits over $K$. Then there exists an embedding of $F$-algebras $K \hookrightarrow A$ of $F$-algebras with (the image of) $K$ a self centralizing maximal subfield of $A$.

**Proof.** By the Remark, there exists a finite dimension central simple $F$-algebra $B$ satisfying $[B] = [A]$ in $\text{Br}(F)$ and $K$ is a self centralizing maximal subfield of $B$. Therefore, $[B : F] = [K : F]^2 = [A : F]$. It follows that $A \cong B$ as $F$-algebras. \hfill \box

**Corollary 104.5.** Let $A$ be a finite dimensional central simple $F$-algebra. Then there exists a finite Galois extension $L/F$ such that $L$ splits $A$ and $L$ is a self centralizing maximal subfield of a finite dimensional central simple $F$-algebra $B$ with $B$ Brauer equivalent to $A$.

**Proof.** We know that there exists a finite Galois extension $L/F$ such that $A$ splits over $L$. Consequently, $L$ splits $D$ if $A \cong M_s(D)$, $D$ a division $F$-algebra. The Remark shows that $L$ is a self centralizing maximal subfield of a finite dimensional central simple $F$-algebra $B$ satisfying $B$ is Brauer equivalent to $D$ in $\text{Br}(F)$. The result follows. \hfill \box

We now investigate a group isomorphic to the Brauer group. We begin by generalizing our construction of the cyclic $F$-algebras in Section 102 that we mentioned in that section. We use the following
Notation 104.6. If a (multiplicative) group acts on a field $L$, we shall denote $\sigma(x)$ by $^\sigma x$ for $x \in L$ and $\sigma \in G$ when convenient.

Let $L/F$ be a finite Galois extension. We now look (again) at the $F$-algebra called a crossed product algebra $A$ mentioned in Section 102. We leave many details left to the reader.

Construction 104.7. Let $L/F$ be a finite Galois extension and $A$ be a vector space over $L$ of dimension $[L : F]$ on basis $B := \{u_\sigma \mid \sigma \in G(L/F)\}$. So $A = \bigoplus_{\sigma \in G(L/F)} u_\sigma$. We define multiplication on $A$ as follows: Let $\sigma, \tau, \eta$ be arbitrary elements in $G(L/F)$ and set
1. $u_\sigma x = \sigma(x)u_\sigma = ^\sigma xu_\sigma$ for all $x \in L$ defines an $L$-action on $A$.
2. $u_\sigma u_\tau = f(\sigma, \tau)u_{\sigma\tau}$ for some $f(\sigma, \tau) \in L^\times$.

In particular, $f : G(L/F) \times G(L/F) \to L^\times$ and $f$ commutes with all elements in $F$. We define multiplication in $A$ by

$$\left( \sum_{\sigma \in G(L/F)} x_\sigma u_\sigma \right) \left( \sum_{\tau \in G(L/F)} y_\tau u_\tau \right) = \sum_{\sigma, \tau \in G(L/F)} x_\sigma y_\tau f(\sigma, \tau)u_{\sigma\tau}.$$

In order for $A$ to be associative, we must have

$$u_\rho(u_\sigma u_\tau) = (u_\rho u_\sigma)u_\tau$$

for all $\rho, \sigma, \tau \in G(L/F)$. As

$$u_\rho(u_\sigma u_\tau) = u_\rho f(\sigma, \tau)u_{\sigma\tau} = ^\rho f(\sigma, \tau)u_{\rho\sigma\tau} = f(\rho, \sigma\tau)u_{\rho\sigma\tau}.$$

and

$$(u_\rho u_\sigma)u_\tau = f(\rho, \sigma)u_{\rho\sigma}u_\tau = f(\rho, \sigma)f(\rho\sigma, \tau)u_{\rho\sigma\tau}$$

for all $\rho, \sigma, \tau \in G(L/F)$, the map $f$ must satisfy

(104.8) $^\rho f(\sigma, \tau)f(\rho, \sigma\tau) = f(\rho, \sigma)f(\rho\sigma, \tau)$

for all $\rho, \sigma, \tau \in G(L/F)$. Check that under these conditions that $A$ is an $F$-algebra.

Definition 104.9. Let $L/F$ be a finite Galois extension. Call a map $f : G(L/F) \times G(L/F) \to L^\times$ satisfying Equation 104.8 a factor set. The algebra $A$ in the construction above is called a crossed product algebra on $f$, and is denoted $(L/F, f)$. We also call $\{u_\sigma \mid \sigma \in G(L/F)\}$ a canonical basis for $(L/F, f)$.

Example 104.10. Every cyclic $F$-algebra is a crossed product algebra. [Cf. the special case in Example 102.4(3).]
Construction 104.11. Let \( f, g := G(L/F) \times G(L/F) \to L^\times \) with \( L/F \) a finite Galois extension. Define \( f, g : G(L/F) \times G(L/F) \to L^\times \) by
\[
(fg)(\sigma, \tau) := f(\sigma, \tau)g(\sigma, \tau) \quad \text{for all } \sigma, \tau \in G(L/F).
\]
The set
\[
Z^2(G(L/F), L^\times) := \{ f : G(L/F) \times G(L/F) \to L^\times \mid f \text{ a factor set} \}
\]
is an abelian group with identity the trivial factor set \( 1 \) defined by \( 1(\sigma, \tau) = 1 \) for all \( \sigma, \tau \in G(L/K) \).

Let \( c : G(K/F) \to L^\times \). Define
\[
\delta c : G(K/F) \times G(K/F) \to L^\times \text{ by } (\delta c)(\sigma, \tau) := c(\sigma)c(\tau)c(\sigma \tau)^{-1}.
\]

Check. \( \delta c \) lies in \( Z^2(G(L/F), L^\times) \).

It is called a Principal Factor Set. Set
\[
B^2(G(L/F), L^\times) := \{ \delta c \mid c : G(L/F) \to L^\times \},
\]
a subgroup of \( Z^2(G(L/F), L^\times) \).

The quotient group
\[
H^2(G(L/F), L^\times) := Z^2(G(L/F), L^\times)/B^2(G(L/F), L^\times)
\]
is called the 2nd cohomology group of \( G(L/F) \) with coefficients in \( L^\times \).

Denote the class of \( f \in Z^2(G(L/F), L^\times) \) in this quotient by \([f]\). If \([f] = [g]\) with \( f, g \in Z^2(G(L/F), L^\times)\), write \( f \sim g \).

Check. If \( f \sim g \), then \((L/F, f) \to (L/F, ((\delta c)f)\) given by \( c(\sigma)u_\sigma \mapsto v_\sigma \), where \( \{u_\sigma \mid \sigma \in G(L/F)\} \) and \( \{v_\sigma \mid \sigma \in G(L/F)\} \) are invariant bases for \((L/F, f)\) and \((L/F, (\delta c)f)\), respectively, is an \( F \)-algebra isomorphism.

Let \( L/F \) be a finite Galois extension. We say \( f \in Z^2(G(L/F), L^\times) \) is a normalized factor set if \( f(\sigma, 1) = 1 = f(1, \sigma) \) for all \( \sigma \in G(L/F) \).

Lemma 104.12. Let \( L/F \) be a finite Galois extension and \( f \) an element in \( Z^2(G(L/F), L^\times) \). Then there exists a normalized factor set \( g \in Z^2(G(L/F), L^\times) \) satisfying \( f \sim g \). In particular, \((L/F, f) \cong (L/F, g)\) as \( F \)-algebras.

Proof. Define \( c : G(L/F) \to L^\times \) by
\[
c(\sigma) = \begin{cases} f(1, 1)^{-1} & \text{if } \sigma = 1 \\ 1 & \text{if } \sigma \neq 1. \end{cases}
\]

Then \((\delta c)(\sigma, 1) = c(\sigma) \) for all \( \sigma \in G(L/F) \). Check if \( g = (\delta c)f \), then \( g \) is a normalized factor set. \( \square \)
Proposition 104.13. Let \( A = (L/K, f) \) with \( f \) a normalized factor set, \( \{ u_\sigma \mid \sigma \in G(L/F) \} \) a canonical basis. Then we have \( u_\sigma \in A^\times \) for all \( \sigma \in G(L/F) \) and identifying \( L \) and \( Lu_1 \) that \( u_1 = 1_A \). [If \( f \) is not normalized, we still have \( u_\sigma \in A^\times \) for all \( \sigma \in G(L/F) \).]

**Proof.** Let \( \{ u_\sigma \mid \sigma \in G(L/F) \} \) be a canonical basis for \( f \), so \( u_1 = 1 \). We know that \( F \subset Z(A) \), i.e., \( A \) is an \( F \)-algebra. Let \( z = \sum_{\sigma \in G(L/F)} a_\sigma u_\sigma \in Z_A(L) \) with all \( a_\sigma \in L \). Then for all \( x \in L \), we have

\[
\sum_{\sigma \in G(L/F)} x a_\sigma u_\sigma = xz = zx = \sum_{\sigma \in G(L/F)} a_\sigma u_\sigma x = \sum_{\sigma \in G(L/F)} a_\sigma \sigma(x) u_\sigma.
\]

Therefore, \( xa_\sigma = \sigma(x) a_\sigma \) for all \( \sigma \in G(L/F) \) and all \( x \in L \). Since \( F = L^{G(L/F)} \), we must have \( a_\sigma = 0 \) for all \( \sigma \neq 1 \) in \( L \). Hence \( z = a_1 u_1 \in L \), so \( L = Z_A(L) \) and we must have \( L \subset A \) is a maximal subfield. Moreover, \( [A : L] = |G(L/F)| = [L : F] \), so \( [A : L] = [L : F]^2 \). In particular, if we show \( A \) is a central simple \( F \)-algebra, then \( L \) is a self centralizing maximal subfield.

We first show that \( A \) is central. Let \( z \in Z(A) \subset Z_A(L) = L \). Then \( zu_\sigma = u_\sigma z = \sigma(z) u_\sigma \) for all \( \sigma \in G(L/F) \). As \( u_\sigma \in A^\times \) for all \( \sigma \in G(L/F) \), we have \( z = \sigma(z) \) for all \( \sigma \in G(L/F) \). Therefore, \( z \in F \), i.e., \( F = Z(A) \), and \( A \) is \( F \)-central. Finally, we show that \( A \) is simple. If \( \mathfrak{A} \) is a nonzero ideal in \( A \), choose 0 \( \neq s = \sum_{i=1}^r a_\sigma_i u_\sigma_i \) in \( \mathfrak{A} \) with \( r \) minimal. Suppose that \( r > 1 \). Choose \( b \in L \) to satisfy \( \sigma_1(b) \neq \sigma_2(b) \). Then we have

\[
0 \neq s - \sigma_1(b)^{-1} s b = \sum_{i=2}^r (a_{\sigma_i} - \sigma_1(b)^{-1} a_{\sigma_i} \sigma_i(b)) u_{\sigma_i}
\]

lies in \( \mathfrak{A} \), contradicting the choice of \( r \). Consequently, \( r = 1 \) and \( 0 \neq s = a_{\sigma_1} u_{\sigma_1} \) lies in \( A^\times \). Therefore, \( \mathfrak{A} = A \). In particular, \( A \) is simple. \( \square \)

**Corollary 104.14.** \( L/F \) be a finite Galois extension and \( f \) an element in \( Z^2(G(L/F), L^\times) \). Let \( A = (L/F, f) \). Then \( A \) is a finite dimensional central simple \( F \)-algebra and \( L \) is a centralizing maximal subfield of \( A \).

**Proof.** \( A \) is isomorphic to a crossed product algebra defined by a normalized factor set. \( \square \)
Proposition 104.15. Let $L/F$ be a finite Galois extension and $f, g \in Z^2(G(L/F), L^\times)$. Then $[f] = [g]$ in $H^2(G(L/F), L^\times)$ if and only if $(L/F, f) \cong (L/F, g)$ as $F$-algebras.

Proof. ($\Rightarrow$) has been done.

($\Leftarrow$): We may assume that both $f$ and $g$ are normalized factor sets with $\{u_\sigma \mid \sigma \in G(L/F)\}$ a canonical basis for $A := (L/F, f)$ and $\{v_\sigma \mid \sigma \in G(L/F)\}$ a canonical basis for $B := (L/F, g)$. In particular, $w_1 = 1_A$ and $v_1 = 1_B$. Let $\varphi : A \to B$ be an $F$-algebra isomorphism. We must have $\varphi(u_1) = v_1$, so $\varphi(Lu_1) = L'v_i$ with $L' \cong L$ as $F$-algebras.

Let $u'_\sigma = \varphi(u_\sigma)$ for all $\sigma \in G(L/K)$, so $\varphi(Lu_\sigma) = L'u'_\sigma$. Therefore, we have $B = \coprod_{G(L/F)} Lu'_\sigma$ and equations

$$
\begin{align*}
u'_\sigma \varphi(x) &= \varphi(c(x))u'_\sigma
\end{align*}
$$

(104.16)

$$
\begin{align*}
u'_\sigma u'_\tau &= \varphi(f(\sigma, \tau))u_{\sigma\tau}
\end{align*}
$$

for all $\sigma, \tau \in G(L/F)$. Since $Lv_1 \cong L'v_1$ are both simple subalgebras of $B$, by the Skolem-Noether Theorem 103.18, there exists an inner automorphism $\theta$ of $B$ satisfying $\theta(\varphi(x)) = x$ for all $x \in L$. Define $w'_\sigma = \theta(u'_\sigma)$ for all $\sigma \in G(L/F)$. Then applying the map $\theta$ to the equations 104.16 yields

$$
B = \coprod_{G(L/F)} L w'_\sigma
$$

$$
\begin{align*}
w_\sigma x &= \sigma(x)w_\sigma
\end{align*}
$$

$$
\begin{align*}
w_\sigma w_\tau &= f(\sigma, \tau)w_{\sigma\tau}
\end{align*}
$$

for all $\sigma, \tau \in G(L/F)$ and $x \in L$. As $v_\sigma v_{\sigma^{-1}} = g(\sigma, \sigma^{-1})v_1 = g(\sigma, \sigma^{-1})$, we have $v_{\sigma^{-1}} = g(\sigma, \sigma^{-1})^{-1}v_{\sigma^{-1}}$ for all $\sigma \in G(L/F)$. Hence for all $x \in L$, we see that

$$
\begin{align*}
w_\sigma v_{\sigma^{-1}} x &= w_\sigma g(\sigma, \sigma^{-1})^{-1}x = w_\sigma g(\sigma, \sigma^{-1})^{-1}\sigma^{-1}(x)v_{\sigma^{-1}}
\end{align*}
$$

$$
\begin{align*}
&= x w_\sigma g(\sigma, \sigma^{-1})^{-1}v_{\sigma^{-1}} = x w_\sigma v_{\sigma^{-1}}.
\end{align*}
$$

It follows that $w_\sigma v_{\sigma^{-1}}$ lies in $Z_B(L) = L$ and $w_\sigma v_{\sigma^{-1}} \neq 0$. Consequently, there exists a map $c : G(L/F) \to L^\times$ satisfying $w_\sigma = c(\sigma)v_\sigma$ for all $\sigma \in G(L/F)$ and

$$
\begin{align*}
f(\sigma, \tau)w_{\sigma\tau} &= w_\sigma w_\tau = c(\sigma)v_\sigma c(\tau)v_\tau = c(\sigma) c(\tau)v_\sigma v_\tau
\end{align*}
$$

$$
\begin{align*}
&= c(\sigma) c(\tau) g(\sigma, \tau) v_{\sigma\tau} = ((\delta c) g)(\sigma, \tau) w_{\sigma\tau}
\end{align*}
$$

for all $\sigma, \tau \in G(L/F)$. Therefore, $f = (\delta c)g$ as needed. □

Corollary 104.17. Let $L/F$ be a finite Galois extension of degree $n$ and $f \in Z^2(G(L/F)/L^\times)$. Then $(L/F, f) \cong \mathbb{M}_n(F)$ if and only if $[f] = 1$ in $H^2(G(L/F), L^\times)$. 

Proposition 104.19. is completely determined by \( [\lambda] \) and \( [\mu] \).

Remark 104.18. If \( A \) is a finite dimensional central simple \( F \)-algebra and \( M \) is a finitely generated \( A \)-module, the \( A \)-isomorphism type of \( M \) is completely determined by \([ M : F ] \). Indeed if \( N \) is an irreducible \( A \)-module, then there exists a unique integer \( n \) such that \( M \cong N^n \) as \( A \)-modules and \([ M : F ] = n[ N : F ] \).

Proposition 104.19. Let \( L/F \) be a finite Galois extension and \( f,g \in Z^2(\mathcal{G}(L/F),L^\times) \). Then \((L/F,f) \otimes_F (L/f,g) \sim (L/F,f,g) \).

Proof. (Chase). Let \( h = fg \). Set

\[
\begin{align*}
A &= (L/K,f) \\
B &= (L/K,g) \\
C &= (L/K,h)
\end{align*}
\]

and \( C = (L/K,h) \) with canonical basis \( \{ v_\sigma \mid \sigma \in \mathcal{G}(L/F) \} \).

Let \( M = A \otimes_L B \) be the tensor product of \( A \) and \( B \) as \( L \)-vector spaces, i.e., \( \alpha a \otimes_L b = \alpha \otimes_L ab \) for all \( \alpha \in L \). We give \( M \) the structure of a right \( A \otimes_F B \)-module by

\[
(a \otimes_L b)(a' \otimes b') := aa' \otimes_L b b' \quad \text{for all } a,a' \in A, b,b' \in B.
\]

Claim. \( M \) is a (left) \( C \)-module by

\[
(cw)_{\sigma}(a \otimes_L b) = (cu_{\sigma}a \otimes_L v_{\sigma}b)
\]

for all \( a \in A, b \in B, c \in C, \) and \( \sigma \in \mathcal{G}(L/F) \):

The only thing that is not straightforward is associativity, i.e., \( (\lambda \lambda')m = \lambda(\lambda' m) \) for all \( \lambda, \lambda' \in C, m \in M \) which we show, leaving the rest as an exercise. To show this, it suffice to assume that \( \lambda = cw, \lambda' = c'w', \) and \( m = a \otimes_L b, \) for \( c,c' \in L, \sigma, \tau \in \mathcal{G}(L/F), a \in A, \) and \( b \in B. \) But

\[
(\lambda \lambda')m = c'c'hw(\sigma,\tau)v_{\sigma\tau}(a \otimes_L b)
\]

\[
= c'c'f(\sigma,\tau)u_{\sigma\tau}a \otimes_L g(\sigma,\tau)v_{\sigma\tau}b = c'c'u_{\sigma}u_{\tau}a \otimes_L v_{\sigma}v_{\tau}b
\]

\[
= \lambda(c'u_{\tau}a \otimes_L v_{\sigma}v_{\tau}b) = \lambda(\lambda' m)
\]

as needed, and the Claim is established.

It follows that \( M \) is a \( C - (A \otimes_F B) \)-bimodule. This means that the map \( \rho : A \otimes_F B \to \mathcal{E}nd_C(M) \) by \( a \otimes b \mapsto \rho_a \otimes b, \) i.e., right multiplication by \( a \otimes b, \) for all \( a \in A, b \in B \) is an \( F \)-algebra homomorphism.
Let $n = [L : F] = [A : L] = [B : L][C : L]$. Then
$$[M : F] = [M : L][L : F] = n^3 = n[C : F].$$
Since $C$ is a finite dimensional central simple $F$-algebra, $M \cong C^n$ as $C$-modules by Remark 104.18, hence $\text{End}_C(M) \cong M_n(C)$. It follows that $[M_n(C) : F] = n^2[C : F] = n^4 = [A \otimes F B : F]$. Therefore, the $F$-algebra homomorphism $\rho$ must be an isomorphism. 

**Theorem 104.20.** Let $L/F$ be a finite Galois extension. Then
$$\Phi_{L/F} : H^2(G(L/F), L^\times) \rightarrow \text{Br}(L/F)$$
defined by $[f] \mapsto (L/F, f)$ is a group isomorphism.

**Proof.** By Proposition 104.19, $\Phi$ is a group homomorphism, hence by Corollary 104.17, it injective. So we need only show that it is onto. Let $A$ be a finite dimensional central simple $F$-algebra split by $L$. Replacing $A$ by a (Brauer) equivalent algebra, we may assume that $L \subset A$ is a simple subalgebra. By the Skolem-Noether Theorem 103.18, for each $\sigma \in G(L/F)$, there exists an element $u_\sigma \in A^\times$ satisfying $\sigma(x) = u_\sigma x u_\sigma^{-1}$ for all $x \in L$, i.e., we have
$$u_\sigma = \sigma x u_\sigma$$
for all $\sigma \in G$, all $x \in L$.

**Claim.** $A = \bigsqcup_{G(L/F)} Lu_\sigma$ as an $L$-vector space:

If we show $\bigsqcup_{G(L/F)} Lu_\sigma = \sum_{G(L/F)} Lu_\sigma$, then this is true, since we know that $[A : F] = [L : F]^2$ and $\bigsqcup_{G(L/F)} Lu_\sigma \subset A$. Suppose that $B = \{u_\sigma \mid \sigma \in G(L/F)\}$ is $L$-linearly dependent. Let $w = \sum_{i=1}^r a_{\sigma_i} u_{\sigma_i} = 0$ with $u_{\sigma_i} \in B$, $a_{\sigma_i} \in L$ be chosen with $r$ minimal. Each $u_\sigma \in A^\times$ for $\sigma \in G(L/F)$, $r > 1$. Choose $b \in L$ satisfying $\sigma_1(b) \neq \sigma_2(b)$. Then
$$0 = \sigma_1(b) w - w b = \sum_{i=2}^r \left(\sigma_1(b) - \sigma_i(b)\right) a_{\sigma_i} u_{\sigma_i},$$
contracts the minimality of $r$. This establishes the Claim.

Equation (*) implies that $u_\sigma u_\tau u_{\sigma \tau}^{-1}$ lies in $Z_A(L) = L$. Hence there exist $f(\sigma, \tau)$ in $L^\times$ satisfying $u_\sigma u_\tau = f(\sigma, \tau) u_{\sigma \tau}$ for all $\sigma, \tau \in G(L/K)$. The associativity for $A$ now implies that $f \in Z^2(G(L/F), L^\times)$. 

**Remark 104.21.** Amitsur has shown that there exist finite dimensional central division $F$-algebras that are not isomorphic to cross product algebras.

**Proposition 104.22.** Let $A$ be a finite dimensional central simple $F$-algebra of index $n$. Then $[A^n] = 1$ in $\text{Br}(F)$. 


Proof. By Theorem 104.20 and Corollary 104.14, we may assume that \( A = (L/F, f) \) with \( L/F \) a finite Galois extension say of degree \( m \) and \( L \) a self centralizing maximal subfield of \( A \) with canonical basis \( \{ u_\sigma \mid \sigma \in G(L/F) \} \) for \( A \). In particular, we have \( u_\sigma x = \sigma(x)u_\sigma \) for all \( x \in L \). We have \( A \cong \mathcal{M}_m(D) \) for some central division \( F \)-algebra and \( \text{ind} A \mid \text{deg} A \). We also have \( \text{deg} A = [L : F] = mn \). Let \( M \) be a simple \( A \)-module, so \( M \) is a \( D \)-vector of dimension \( m \) (as isomorphic to a column of \( \mathcal{M}_m(D) \)). Then

\[
mn \dim_L M = [L : F] \dim_L M = \dim_F M = \dim_D M \dim_F D = mn^2,
\]

so \( \dim_L M = n \). Let \( \{ v_1, \ldots, v_n \} \) be an \( L \)-basis for \( M \). For each \( x \in A \), define \( C(x) \in \mathcal{M}_m(D) \) by

\[
xv_j = \sum_{i=1}^{m} C(x)_{ij} v_i.
\]

We have

(1) \( u_\sigma u_\tau v_j = f(\sigma, \tau) u_{\sigma \tau} v_j = f(\sigma, \tau) \sum_{i=1}^{m} C(u_{\sigma \tau})_{ij} v_i \)

and

\[
u_\sigma u_\tau v_j = u_\sigma \left( \sum_{k=1}^{m} C(u_\tau)_{kj} v_k \right) = \left( \sum_{k=1}^{m} \sigma(C(u_\tau)_{kj}) u_\sigma v_k \right)
= \sum_{i,k=1}^{m} \sigma(C(u_\tau)_{kj}) (C(u_\sigma)_{ik} v_i).
\]

Let \( \sigma(C) = (\sigma(C_{ij})) \) in \( \mathcal{M}_m(F) \), i.e., the matrix in which \( \sigma \) acts on every entry of \( C \). So we have

(2) \( u_\sigma u_\tau v_j = \sum_{i=1}^{m} \sigma(C(u_\tau)_{ij}) v_i \).

Comparing (1) and (2) yields

\[
f(\sigma, \tau)C(u_{\sigma \tau}) = C(u_\sigma)\sigma(C(u_\tau))
\]

in \( \mathcal{M}_m(L) \). Setting \( c_\sigma = \det(C(u_\sigma)) \) for all \( \sigma \in G(L/F) \), we deduce that

\[
f(\sigma, \tau)^m c_{\sigma \tau} = c_\sigma \sigma(c_\tau)
\]

for all \( \sigma, \tau \in G(L/F) \). It follows that \( f(\sigma, \tau)^m = c_\sigma \sigma(c_\tau) c_{\sigma \tau}^{-1} \) is the trivial factor set.

\[\square\]

Corollary 104.23. The Brauer group of a field is a torsion group.
**Proof.** \( \text{Br}(F) = \bigcup_{L/F \text{ finite Galois}} \text{Br}(L/F) \) \( \square \)

**Definition 104.24.** Let \( A \) be a a finite dimensional central simple \( F \)-algebra. The order of \([A]\) in the Brauer group of \( F \) is called the *exponent* or *period* of \( A \) and is denoted \( \text{exp} \ A \).

**Corollary 104.25.** Let \( A \) be a a finite dimensional central simple \( F \)-algebra. Then \( \text{exp} \ A \), the exponent of \( A \), divides \( \text{ind} \ A \), the index of \( A \).

**Remark 104.26.** If \( A \) is a a finite dimensional central simple \( F \)-algebra, then the \( \text{exp} \ A \) and \( \text{ind} \ A \) have the same prime divisors. We leave this as an exercise.

**Remark 104.27.** We indicate (with few details) how one can put together all the \( H^2(G(L/F), L^\times) \) with \( L/F \) a finite Galois extension to get a 2nd cohomology group isomorphic to the Brauer group \( \text{Br}(F) \).

Let \( \tilde{F} \) be an algebraic closure of \( F \) and \( \Gamma = G(F_{\text{sep}}/F) \), the (absolute) Galois group of \( F \), where \( F_{\text{sep}} \) is the separable closure of \( F \) in \( \tilde{F} \). We saw that \( \Gamma \) has a topology, the profinite topology. Let \( F^\times_{\text{sep}} \) have the discrete topology, and define

\[
\begin{align*}
Z^2_{\text{cont}}(\Gamma_F, F^\times_{\text{sep}}) &:= \{ f \in Z^2(\Gamma_F, F^\times_{\text{sep}}) \mid f \text{ continuous} \}, \\
B^2_{\text{cont}}(\Gamma_F, F^\times_{\text{sep}}) &:= \{ f \in B^2(\Gamma_F, F^\times_{\text{sep}}) \mid f \text{ continuous} \}, \\
H^2_{\text{cont}}(\Gamma_F, F^\times_{\text{sep}}) &:= Z^2_{\text{cont}}(\Gamma_F, F^\times_{\text{sep}})/B^2_{\text{cont}}(\Gamma_F, F^\times_{\text{sep}}).
\end{align*}
\]

Let \( E/L/F \) be finite extensions in \( \Gamma_F \) with \( L/F \) and \( E/F \) Galois extensions. One shows that the group inclusion \( i_{E/F} : \text{Br}(L/F) \to \text{Br}(E/F) \) and the group epimorphism \( G(E/F) \to G(L/F) \) given by \( \sigma \mapsto \sigma|_L \) (so \( G(E/F)/G(E/L) \cong G(L/F) \)) of Galois theory induce a group monomorphism

\[
\inf_{E/L} : H^2(G(L/F), L^\times) \to H^2(G(E/F), E^\times)
\]

such that

\[
\begin{array}{ccc}
\text{Br}(L/F) & \xrightarrow{i_{E/L}} & \text{Br}(E/F) \\
\Phi_{L/F} & & \Phi_{E/F} \\
H^2(G(L/F), L^\times) & \xrightarrow{\inf_{E/L}} & H^2(G(E/F), E^\times).
\end{array}
\]

commutes. One then can show that \( H^2_{\text{cont}}(\Gamma_F, F^\times_{\text{sep}}) \) is built from the \( H^2(G(L/F), L^\times) \) as \( L \) runs over the finite Galois extensions \( L/F \). More
specifically, \( H^2_{cont}(\Gamma_F, F^\times_{sep}) = \lim_\to H^2(G(L/F), L^\times) \), the direct limit over \( \inf_{L/F} \) where the \( L/F \) run over finite Galois extensions in \( \tilde{F} \). (Cf. Exercise 104.30(7) for the definition of direct limit.) We then obtain a group isomorphism \( \Phi_F : Br(F) \to H^2_{cont}(\Gamma_F, F^\times_{sep}) \) compatible with all \( \Phi_{L/F} \).

We end this section with remarks amplifying comments made about cyclic algebras in Section 102. We omit the proofs, since we shall not develop cohomology theory.

**Remark 104.28.** Let \( L/F \) be a cyclic field extension of degree \( n \) with \( G(L/F) = \langle \sigma \rangle \) and \( A = (L/F, \sigma, a) \) a cyclic \( F \)-algebra. Then the following can be shown:

1. Suppose that \( L/K/F \) is an intermediate field with \( [K : F] = m \). Then \( (L/F, \sigma|_K, a) \cong (L/K, \sigma, a^{\frac{1}{m}}) \). In particular, \( A \) is a division \( F \)-algebra if \( \pi \) has order \( n \) in \( F^\times/N_{L/F}(L^\times) \). The last statement follows as we have \( n = \exp A = \deg A \) in this case (cf. the exercises for this section).

2. There is a group isomorphism \( F^\times/N_{L/F}(L^\times) \to H^2(G(L/F), L^\times) \) mapping \( xN_{L/F}(L^\times) \mapsto (L/F, \sigma, x) \). [This arises from the periodicity of Tate cohomology for finite cyclic groups.]

3. \( A \) is a central division \( F \)-algebra if \( [A] \) has order \( n \) in \( Br(L/F) \).

**Remark 104.29.** Historically, much of our study of finite dimensional central simple algebras reached a high point in algebraic number theory, where the classification of finite dimensional central division \( F \)-algebras with \( F/\mathbb{Q} \) is a finite extension was determined in a joint paper of Brauer, Hasse, and Noether (some of the results were independently proven by Albert). Brauer had defined cross product algebras (and the Brauer group) which was simplified by Noether, generalizing cyclic algebras over arbitrary fields. Recall if \( F/\mathbb{Q} \) is a finite extension (an algebraic number field), then

\[
\mathbb{Z}_F := \{ x \in K \mid f(x) = 0, \text{ for some monic polynomial } f \in \mathbb{Z}[t] \}
\]

is the ring of algebraic integers in \( F \). If \( p \) is a prime ideal in \( \mathbb{Z}_F \), let \( (\mathbb{Z}_F)_p \) denote its completion (cf. Exercise 95.32(4)) at \( p \) and \( F_p \) the quotient field of \( (\mathbb{Z}_F)_p \). All central simple \( F_p \)-algebras were known to be cyclic, in fact, one central division ring for each integer \( n \) (up to isomorphism). Let \( K/F \) be a finite extension. Then \( p\mathbb{Z}_K \) factors into a product of primes in \( \mathbb{Z}_K \). If \( \mathfrak{p} \) is one such, we write \( p \mid \mathfrak{p} \) and we get a field extension \( K_{\mathfrak{p}}/F_p \). Let \( X(F) = \text{Max}\mathbb{Z}_K \cup X_\infty(\mathbb{Z}_K) \) where \( X_\infty(\mathbb{Z}_K) := \{ f : \mathbb{Z}_F \to \mathbb{C} \mid f \text{ a field homomorphism} \} \). So each
\( f \in X_\infty(Z_K) \) has image in \( \mathbb{R} \), with completion \( \mathbb{R} \), called a real infinite prime or not, called a complex infinite prime, with completion \( \mathbb{C} \). We then extend our notation to all elements in \( X(F) \). Let \( L/F \) be a cyclic extension. Hasse proved that \( x \in F \) satisfies \( x \in N_{L/F}(L) \) if and only if \( x \in N_{L_\mathfrak{p}/F_\mathfrak{p}}(L_\mathfrak{p}) \) for all \( \mathfrak{p} \in X(F) \) and all \( \mathfrak{p} \in X(L) \) with \( \mathfrak{p} \mid \mathfrak{p} \).

One says that \( x \) is a norm if and only if \( x \) is a norm locally for all primes (finite and infinite). As cyclic \( F \)-algebras split under a norm condition, this local-global principle applies to cyclic \( F \)-algebras. They then showed every crossed product \( F \)-algebra was a cyclic \( F \)-algebra and every central simple \( F \)-algebra was a cyclic algebra.

**Exercises 104.30.**

1. Let \( A \) and \( B \) be finite dimensional central simple \( F \)-algebras. Show the following:
   (a) Let \( K/F \) is a finite field extension with \( \text{ind} \ A \) relatively prime to \( [K : F] \). Then \( \text{ind} \ A^K \mid \text{ind} \ A \). In particular, if in addition, \( A \) is a division algebra, so is \( A^K \).
   (b) We have \( \text{ind}(A \otimes_F B) \) divides \( (\text{ind} \ A)(\text{ind} \ B) \).
   (c) We have \( \text{ind}(A \otimes_F \cdots \otimes_F A)_m^{\text{m divides ind} A} \) divides \( \text{ind} A \) for all \( m \geq 1 \).

2. Let \( D \) be a central division \( F \)-algebra and \( K/F \) a finite field extension with \( [K : F] \) a prime dividing \( \text{deg} \ D \). Prove that the following are equivalent:
   (a) \( K \) is isomorphic to a subfield of \( D \).
   (b) \( D^K \) is not a division algebra.
   (c) \( \text{deg} \ D = [K : F] \text{ind} D^K \).

3. Let \( A \) be a finite dimensional central simple \( F \)-algebra with \( \text{ind} \ A = p^m e \), \( p \) a prime not dividing \( m \) and \( e \geq 1 \). Prove that there exists a finite separable extension \( K/F \) satisfying \( p \nmid [K : F] \) and \( \text{ind} A^K = p^e \).

4. Show Example 104.10 is a cross product algebra.

5. Let \( A \) and \( B \) be finite dimensional central simple \( F \)-algebras and \( K/F \) a finite field extension. Show the following:
   (a) Every prime divisor of \( \text{ind} \ A \) divides \( \text{exp} \ A \).
   (b) We have \( \text{exp} A^K \mid \text{exp} A \).
   (c) If \( \text{ind} A \) and \( [K : F] \) are relatively prime, then \( \text{exp} A^K = \text{exp} A \).
   (d) We have \( \text{exp}(A \otimes_F B) \) divides the least common multiple of \( \text{exp} A \) and \( \text{exp} B \).
   (e) We have \( \text{exp}(A \otimes_F \cdots \otimes_F A)_m^{\text{m divides exp} A} = (\text{exp} A)/d \) where \( d \) is the gcd of \( \text{exp} A \) and \( [K : F] \).
(f) If \( \text{ind} A \) and \( \text{ind} B \) are relatively prime, then \( (\text{ind}(A \otimes_F B) = (\text{ind} A)(\text{ind} B) \) and \( (\text{exp}(A \otimes_F B) = (\text{exp} A)(\text{exp} B) \). In particular, if in addition, \( A \) and \( B \) are division algebras, so is \( A \otimes_F B \).

6. Let \( D \) be a central division \( F \)-algebra with \( \text{deg} D = p_1^{e_1} \cdots p_r^{e_r} \) its standard factorization. Then there exists an \( F \)-algebra isomorphism \( D = D_1 \otimes_F \cdots \otimes_F D_r \) with \( D_i \) a central \( F \)-division algebra with \( \text{deg} D_i = p_i^{e_i} \) for all \( i = 1, \ldots, r \), unique up to isomorphism.

7. Let \( I \) be a partially ordered set under \( \leq \) that also satisfies for all \( i, j \in I \), there exists a \( k \in I \) such that \( i \leq k \) and \( j \leq k \); and \( \{M_i\}_I \) be a collection of \( R \)-modules. Suppose for all \( i \leq j \) in \( I \), there exist \( R \)-homomorphisms \( \theta_{i,j} : M_i \to M_j \) with \( \theta_{i,i} = 1_{M_i} \) and

\[
\begin{array}{ccc}
M_i & \xrightarrow{\theta_{i,k}} & M_k \\
\downarrow{\theta_{ij}} & & \downarrow{\theta_{jk}} \\
M_j & \xrightarrow{\theta_{jk}} & M_k
\end{array}
\]

commutes whenever \( i \leq j \leq k \). Show that there exists an \( R \)-module \( M \) and for each \( i \in I \) an \( R \)-homomorphism \( \psi_i : M_i \to M \) satisfying the following universal property: For all non-negative \( i \leq j \) in \( I \),

\[
\begin{array}{ccc}
M_i & \xrightarrow{\psi_i} & M \\
\downarrow{\theta_{ij}} & & \downarrow{\psi_j} \\
M_j & \xrightarrow{\psi_j} & M
\end{array}
\]

commutes, and if there exist \( R \)-homomorphisms \( \theta_i : M_i \to M \) satisfying for all \( i \leq j \) in \( I \),

\[
\begin{array}{ccc}
M_i & \xrightarrow{\varphi_i} & M' \\
\downarrow{\theta_{ij}} & & \downarrow{\varphi_j} \\
M_j & \xrightarrow{\varphi_j} & M'
\end{array}
\]
commutes, then there exists a unique $R$-homomorphism $\mu: M_i \to M'$ satisfying

$$M_i \xrightarrow{\theta_{ij}} M \xrightarrow{\phi_i} M_j \xrightarrow{\psi_j} M'$$

commutes for all $i \leq j$ in $I$. Such a $M$ is unique up to a unique isomorphism and is called the direct limit of the $M_i$ and denoted by $\lim \rightarrow M_i$.

8. Show reversing all the arrows in Exercise 7 still gives a valid result and the unique $R$-module $M$ up to isomorphism is called the inverse or projective limit of the $M_i$ and denoted by $\lim \leftarrow M_i$.

[Remark: For example, if $\Gamma_F$ is the absolute Galois group of a field $F$ (in some fixed algebraic closure), then $\Gamma_F = \lim \leftarrow G(L/F)$, where the inverse limit is over all finite Galois extensions $L/F$. We know that $G(L/F) = \Gamma_F/\Gamma_L$ for all Galois extension $L/F$, and we have

$$H^2_{cont}(\Gamma_F, F_{\text{sep}}^\times) = \lim \rightarrow H^2(G(L/F), L^\times) = H^2(\lim \leftarrow G(L/F), L^\times)$$

where $L/F$ runs over all finite Galois extensions of $F$.]

105. Polynomial Rings over a Division Algebra

When we investigated polynomial rings over a commutative ring, we noted that the analogous theory broke down if the ring was not commutative. As an example, we showed that the polynomial $t^2 - 1$ over the division ring of Hamiltonian quaternions had infinitely many roots. (Cf. Remark 33.10(2).) A major reason for this is that if $R$ is a ring that is not commutative and $x$ an element in $R$, then the evaluation map $e_x : R[t] \to R$ given by $\sum_{i=0}^n a_i t^i \mapsto \sum_{i=0}^n a_i x^i$ is not necessarily a ring homomorphism, although it is an additive homomorphism. For example, if with $ab \neq ba$ and $f = (t - a)(t - b)$ in $R[t]$, then

$$f = t^2 - (a + b)t + ab, \text{ hence } f(a) = -ba + ab \neq 0,$$

but $e_a(t - a)e_a(t - b) = 0$. The problem arises as $t$ is central in $R[t]$, i.e., $at = ta$ for all $a \in R$. One useful result, still holding in $R[t]$, is the division algorithm as the usual proof works, i.e., if $f \in R[t]$ is monic (or leading term a unit), for any $g \in R[t]$, there exists $h, r \in R[t]$ satisfying

$$g = hf + r \text{ with } r = 0 \text{ or } \deg r < \deg f.$$
As mentioned in the special case in Appendix C, Lemma C.5, there are in fact two such algorithms, i.e., left and right, so there exist \( h', r' \) in \( R[t] \) with \( g = f h' + r' \) with \( r' = 0 \) or \( \deg r' < \deg f \). We shall concentrate on the first, which we call \textit{right division} by \( f \). As in Appendix C, the special case when \( f \) is linear, say \( f = t - x \), is of interest, i.e.,

\[
g = h \cdot (t - x) + r \quad \text{with} \quad r \in R.
\]

We call \( x \) a \textit{right root} of \( g \in R[t] \) if \( g(x) = 0 \) in the above, i.e., if \( g = \sum_{i=0}^{n} a_i t^i \), then \( \sum_{i=0}^{n} a_i x^i = 0 \). (So a left root of \( g \) would be \( \sum_{i=0}^{n} a_i x^i = 0 \).) Since we are interested in the map \( e_a : R[t] \to R \) by plugging \( a \) into \( t \) on the right in this section, we shall let \( g(a) := e_a(g) \) to be this evaluation on the right. In particular, if \( x \) is a right root of \( g \), then \( g \in R[t](t - x) \). Our analogy to the commutative case is

**Lemma 105.1.** Let \( R \) be a ring, \( 0 \neq g \in R[t] \). Then \( x \in R \) is a right root of \( f \) if and only if \( t - x \) is a right divisor of \( g \). In particular,

\[ \{ h \in R[t] \mid x \text{ is a right root of } h \} = R[t](t - x). \]

**Proof.** Let

\[
f = \left( \sum_{i=0}^{n} a_i t^i \right) (t - x) = \sum_{i=0}^{n} a_i t^{i+1} - \sum_{i=0}^{n} a_i x t^i.
\]

Then \( x \) is a right root of \( f \), as \( \sum_{i=0}^{n} a_i x^{i+1} - \sum_{i=0}^{n} a_i x x^i = 0 \). \( \square \)

A very useful observation, is the following observation:

**Proposition 105.2.** Let \( D \) be a division ring. Suppose we have an equation \( f = gh \) of polynomials in \( D[t] \) with \( x \in D \) satisfying \( r = h(x) \) is nonzero (so a unit in \( D \)). Then

\[
f(x) = g(r x r^{-1}) h(x).
\]

In particular, if \( x \) is a right root of \( f \) but not of \( h \), then \( r x r^{-1} \) is a right root of \( g \).

**Proof.** Let \( g = \sum_{i=0}^{n} a_i t^i \), hence \( f = \sum_{i=0}^{n} a_i h(t) t^i \). Therefore,

\[
f(x) = \sum_{i=0}^{n} a_i h(x) x^i = \sum_{i=0}^{n} a_i r x^i
\]

\[= \sum_{i=0}^{n} a_i (r x r^{-1})^i r = g(r x r^{-1}) h(x). \quad \square
\]

This results in an analogue of the number of roots for a polynomial over a field, an analogue that often comes up in the study. If \( x \) is a nonzero element in a division ring \( D \), we shall let \( C(x) := \{ r x r^{-1} \mid r \in \} \).
Let $D$ be a division ring and $f$ a polynomial in $D[t]$ of degree $d$. Then the right roots of $f$ lie in at most $n$ conjugacy classes in $D$ (not necessarily distinct). In particular, if $f = (t - a_1) \cdots (t - a_n)$ with $a_1, \ldots, a_n \in D$ and $d$ is a right root of $f$, then $d$ is conjugate to $a_i$ for some $i$, $i = 1, \ldots, n$.

**Proof.** The result is immediate if $n = 1$, so assume that $n \geq 2$. Let $x \in D$ be a right root of $f$. Then $f = g \cdot (t - x)$ with $g \in D[t]$ of degree less than $n$. If $y$ is another root of $f$ different from $x$, then by Proposition 105.2, a conjugate of $y$ is a root of $g$, and the first statement follows by induction. If $g = (t - b_1) \cdots (t - b_{n-1})$, the second statement also follows by Proposition 105.2 and induction. \[\square\]

Just as in field theory, we are interested in the analogue of algebraic elements. If $D$ is a division ring, its center $F$ is a field, and we can look at elements $x$ in $D$ algebraic over $F$. Of course, $F(x)$ is also a subfield of $D$, and if $x$ is algebraic it has an irreducible polynomial. But if $x$ is algebraic over $F$, then every conjugate of $x$ is also algebraic over $F$, since a right root of $g \in D[t]$ implies that $rxr^{-1}$ is a right root of $rg(t)r^{-1}$ for all nonzero $r$ in $D$.

**Definition 105.4.** Let $D$ be a division ring with center $F$. Let $x$ in $D$ and $C$ be the conjugacy class of $x$ in $D$. We say that $C$ is algebraic over $F$ if and only if every element of $C$ is algebraic over $F$. Suppose that $x$ in $D$ is algebraic over $F$ with $m_F(x) \in F[t]$ its minimal polynomial. Then $m_F(x) = m_F(rxr^{-1})$ for every nonzero element $r$ in $D$. We call this polynomial the minimal polynomial of $C$ and denote it by $m_F(C)$.

The analogue of the field case now becomes:

**Lemma 105.5.** Let $D$ be a division ring with center $F$ and $C$ a conjugacy class in $D$ algebraic over $F$. Suppose that $h$ is a nonzero polynomial in $D[t]$ satisfying $h(c) = 0$ for all $c$ in $C$. Then $\deg h \geq \deg m_F(C)$.

**Proof.** Suppose the result is false. Among all counterexamples choose one $h$ in $D[t]$ of minimal degree $m$ with $m < \deg m_F(C)$. Let $h = \sum_{i=0}^{m} a_it^i$. We may assume that $a_m = 1$, i.e., that $h$ is monic. By the commutative case, we may also assume that not all the $a_i$ lie in $F$, say $a_j \notin F$. Choose $b$ in $D$ satisfying $a_jb \neq ba_j$. Let $c \in C$, so by
assumption $\sum_{i=0}^{m} a_i c^i = 0$. Hence
\[
0 = b(\sum_{i=0}^{m} a_i c^i)b^{-1} = \sum_{i=0}^{m} (ba_i b^{-1})(bc^i b^{-1}) \quad \text{with } a_j \neq ba_j b^{-1}.
\]

Then the polynomial $h_1 = \sum_{i=0}^{m-1} (a_i - ba_i b^{-1})t^i$ vanishes on $bC b^{-1} = C$, contradicting the minimality of $m$. \hfill \Box

This lemma implies:

**Proposition 105.6.** Let $D$ be a division ring with center $F$ and $C$ a conjugacy class in $D$ algebraic over $F$. If $h \in D[t]$, then $h$ vanishes on $C$ if and only if $h \in D[t]m_F(C)$.

**Proof.** Set $f = m_F(C)$. By Lemma 105.1, for all $x \in D$, we have $f \in D[t] : (t - x)$. It follows that if $h \in D[t]f$, then $h(x) = 0$ for all $x$ in $C$. Conversely, suppose that $h(x) = 0$ for all $x \in C$. By the (right) division algorithm, $h = qf + g$ in $D[t]$ with $g = 0$ or deg $g < \deg f$. Since $h(x) = 0 = f(x)$ for all $x \in C$, we also have $g(x) = 0$ for all $x \in C$. By Lemma 105.5, we must have $g = 0$, so $h \in D[t]f$. \hfill \Box

**Corollary 105.7.** Let $D$ be an infinite division ring. Then no nonzero polynomial in $D[t]$ vanishes on $D$.

**Proof.** We may assume that $D$ is not a field. Suppose the corollary is false and that $h = t^m + a_{m-1}t^{m-1} + \cdots + a_1 t + a_0$ in $D[t]$ vanishes on all of $D$. In particular, $h(0) = 0$, so $a_0 = 0$. We may also assume that we have chosen $h$ with $m$ minimal. Let $F$ be the center of $D$. By the proof of Lemma 105.5, we see that we must have $a_i \in F$ for all $i$, i.e., $h \in F[t]$. It follows that $F$ must be a finite field. Let $x \in D$. Then $x$ is a root of $h$, so $x$ is algebraic over $F$. In particular, the subfield $F[x]$ of $D$ is a finite extension of $F$, so also a finite field. Hence there exists an integer $n = n(x)$ satisfying $x^n = x$. By Proposition 105.5, $D$ is a field, a contradiction. \hfill \Box

**Lemma 105.8.** Let $D$ be a division ring with center $F$ and $f \in F[t]$. Suppose that there exist polynomials $g_1, g_2 \in D[t]$ satisfying $f = g_1 g_2$. Then $f = g_2 g_1$.

**Proof.** As $f \in F[t] \subset Z(D[t])$ (check), $g_1 g = gg_1 = g_1 g_2 g_1$. Since $D[t]$ is a non-commutative domain (obvious definition), $g = g_2 g_1$. \hfill \Box

We now show that the minimal polynomial of an algebraic conjugacy class in a division ring $D$ splits in $D[t]$.

**Theorem 105.9.** (Wedderburn) Let $D$ be a division ring with center $F$ and $C$ a conjugacy class in $D$ algebraic over $F$. Suppose that
\[ \deg m_F(C) = n. \] Then there exist \( a_1, \ldots, a_n \) in \( D \) satisfying
\[ (*) \quad m_F(C) = (t - a_1) \cdots (t - a_n). \]
Moreover, \( a_1 \in C \) can be arbitrarily chosen and then any cyclic permutation of \((*)\) is still \( m_F(C) \).

**Proof.** Let \( f = m_F(C) \). Fix any element \( a_1 \in C \). Then \( f \in D[t](t - a_1) \) by Lemma 105.1. Write \( f = g \cdot (t - a_r) \cdots (t - a_1) \) in \( D[t] \) with \( a_1, \ldots, a_r \in C \) and \( r \) maximal. Let \( h = (t - a_r) \cdots (t - a_1) \).

We show \( h \) vanishes on \( C \). If not then there exist an \( x \in C \) with \( h(x) \neq 0 \). By Proposition 105.2, there exists an \( a_{r+1} \in C \), a conjugate of \( x \) satisfying \( g(a_{r+1}) = 0 \). It follows that \( f = h \cdot (t - a_{r+1}) \cdots (t - a_1) \)
\( \) for some \( h \in D[t] \) by Lemma 105.1. This contradicts the maximality of \( r \). Thus \( h \) vanishes on \( C \), hence lies in \( D[t]f \) by Proposition 105.6.

It follows that \( r = n \) and \( f = (t - a_n) \cdots (t - a_1) \) in \( D[t] \) by Lemma 105.5. The last statement in the theorem follows by Lemma 105.8. \( \square \)

**Definition 105.10.** A division ring \( D \) is called **right algebraically closed** (respectively, **left algebraically closed**) if every non-constant polynomial in \( D[t] \) has a right (respectively, left) root.

We note that Frobenius’s Theorem for division rings over the reals 101.1 holds over any real closed field using the same proof. In particular, the Hamiltonian quaternions \( H \) are defined over any real closed field \( F \). We use the same notation as in Construction 32.1. In particular, we have the map \( - : H \to H \) given by \( x_0 + x_1 i + x_2 j + x_3 k \mapsto x_0 - x_1 i - x_2 j - x_3 k \) is an antiautomorphism fixing \( F \). This antiautomorphism induces an antisomorphism \( - : H[t] \to H[t] \) given by \( \sum_{i=0}^n a_i t^i \mapsto \sum_{i=0}^n \overline{a_i} t^i \).

In particular, if \( f \in H[t] \), then \( \overline{f} f \) lies in \( F[t] \).

**Proposition 105.11.** Let \( F \) be a real closed field. Then the Hamiltonian quaternions \( H \) over \( F \) is right and left algebraically closed.

**Proof.** Let \( f = \sum_{i=0}^n a_i t^i \) be a nonconstant polynomial in \( H[t] \).

We show that \( f \) has a right root. This is certainly true if \( n = 1 \), so we may assume that \( n > 1 \). As the element \( i \in H \) satisfies, \( F(i) \) is algebraically closed and \( \overline{f} f \) lies in \( F[t] \), it has a (right) root \( \alpha \). If \( \alpha \) is a right root of \( f \), we are done, so we may assume not. Then a conjugate \( \beta \) of \( \alpha \) is a right root of \( \overline{f} \) by Proposition 105.2. As \( \sum_{i=0}^n \overline{\beta} a_i = \sum_{i=0}^n \overline{a_i} \overline{\beta}^i = 0 \), we have \( \overline{\beta} \) is a left root of \( f \). By the left analogue of Proposition 105.2, \( f = (t - \overline{\beta})g \) for some \( g \in D[t] \). Since \( \deg g = n - 1 \), the polynomial \( g \) has a right root in \( H \), hence so does \( f \). Consequently, \( H \) is right algebraically closed. The proof that \( H \) is left algebraically closed is analogous. \( \square \)
Using the Artin-Schreier Theorem 74.2, we can now prove:

**Theorem 105.12.** (Baer) Let $D$ be a non-commutative division ring with center $F$. Suppose that $D$ is finite over $F$, i.e., $\dim_F D < \infty$, and every nonzero polynomial in $F[t]$ has a right root in $D$, e.g., $D$ is right algebraically closed. Then $F$ is real closed and $D = H$.

**Proof.** Let $n = \dim_F D$. Let $f$ be a nonconstant polynomial in $F[t]$. Then $f$ has a right root $\alpha$ in $D$. In particular, $\deg m_F(\alpha) \leq n$ by Lemma 105.5.

**Claim.** $F$ is perfect.

We may assume that $F$ has positive characteristic $p$. Let $L$ be an algebraic closure of $F$ and $\alpha$ be an element of $L$. Let $\alpha^{1/p^n}$ be the unique $p^n$th root of $\alpha$ in $L$. We must show $\alpha^{1/p}$ lies in $F$. Since $F \subset F(\alpha^{1/p}) \subset \cdots \subset F(\alpha^{1/p^n}) \subset \cdots$ and $[F(\alpha^{1/p^n}) : F] \leq n$, there exists an $m$ such that $\alpha^{1/p^{m+1}}$ lies in $F(\alpha^{1/p^n})$, hence $\alpha^{1/p}$ lies in $F(\alpha^{1/p^n})$, as needed.

Now let $E$ be an intermediate field satisfying $L/E/F$ chosen with $E$ a simple extension of $F$ (i.e., is $E = F(x)$ some $x$) of maximal degree. Since $E/F$ is separable, $E$ is determined by a separable polynomial $f$ in $F[t]$ and has a root in $D$. Therefore, $[E : F] \leq n$. In particular, $L/F$ must be a finite extension. Since $D$ is not commutative, $F < L$. By the Artin-Schreier Theorem 74.2, $F$ is real closed and hence $D = H$ by Frobenius’s Theorem 101.1.

**Lemma 105.13.** Let $D$ be a division ring, $V$ a left vector space over $F$, and $M$ a subset of $V$ closed under addition and containing at least two linear independent elements in $V$. Suppose that $\varphi : M \to V$ is an additive map such that for all $x \in M$, $\varphi(x) = \lambda_x x$ for some $\lambda_x \in F$. Then there exists an element $\lambda \in F$ satisfying $\varphi(x) = \lambda x$ for all $x \in M$.

**Proof.** Let $x$ be a nonzero element in $M$. By assumption $M \not\subset Fx$. Let $y$ be any element in $M \setminus (M \cap Fx)$. Then

\[\lambda_x x + \lambda_y y = \varphi(x) + \varphi(y) = \varphi(x+y) = \lambda_{x+y} (x+y) = \lambda_{x+y} x + \lambda_{x+y} y.\]

As $x, y$ are linearly independent, it follows that $\lambda_x = \lambda_{x+y} = \lambda_y$. In particular, if $z$ is a nonzero element of $M$ not in $Fx$, we have $\lambda_x = \lambda_z$. If $z \in Fx \subset M$, then replacing $x$ by $z$ in (*) shows $\lambda_z = \lambda_y = \lambda_x$. The result follows.

**Theorem 105.14.** (Cartan-Brauer-Hua) Let $K \subset D$ be division rings and $M \subset D$ be closed under addition and containing 1. Suppose that $xK \subset Kx$ for all $x$ in $M$. Then either $M \subset K$ or $M \subset Z_D(K)$, the centralizer of $K$ in $D$. In particular, if $xKx^{-1} \subset K$ for all $x$ in $D$, then either $K = D$ or $K$ lies in the center of $D$. 

Proof. Suppose that $M \not\subset K$. Then $M$ contains at least two linearly independent vectors in the right vector space $D$ over $K$. Let $k \in K$. If $x \in M$, then $xk \in xK \subset Kx$, hence $xk = \lambda_x x$ for some $\lambda_x$ in $K$. Applying the lemma to the map $\rho_k : M \to M$ given by $x \mapsto xk$ yields an element $\lambda \in K$ satisfying $xk = \lambda x$ for all $x$ in $M$. Setting $x = 1$ shows $\lambda = k$. It follows that $kx = xk$ for all $x$ in $M$ and $k \in K$. This gives the first statement. Applying the first statement with $M = D$ yields the second statement. \hfill \Box

Theorem 105.15. Let $D$ be a division ring and $x$ an element of $D$ having only finitely many conjugates in $D$. Then $x$ has only one conjugate, i.e., lies in the center of $D$.

Proof. Let $Z_D(x) := Z_{D^x}(x) \cup \{0\} = \{d \in D \mid xd = dx\}$, a division ring lying in $D$. Then the hypothesis means that $[D^x : Z_{D^x}(x)]$ is finite. Each conjugate $axa^{-1}$ of $x$ gives rise to a division algebra $Z_D(axa^{-1})$ with $Z_{D^x}(axa^{-1})$ of finite index in $D$ and there are finitely many such. Let $K$ be the intersection of these finitely many division rings. Then $K$ is a division ring satisfying $[D^x : K^x]$ is finite by Poincaré’s Lemma (Exercise 10.16(6)). Moreover, $K^x \subset D^x$. It follows by the Cartan-Brauer-Hua Theorem above that either $K = D$ or $K \subset Z(D)$. If $x \in Z(D)$, we are done. So suppose that $K \subset Z(D))$. Then we must have $[D^x : Z(D^x)] \leq [D^x : Z_{D^x}(K^x)] < \infty$. In particular, if $Z(D)$ is a finite field, then $D$ is a finite division ring, hence commutative by Proposition 105.5. So we may assume that $Z(D)$ is infinite. Let $y_1, y_2, \ldots$ be an infinite number of elements in $Z(D)$ and $x_0 = x$, $x_1 = x + y_1$, $x_2 = x + y_2$, $\ldots$. Since $[D^x : Z(D^x)]$ is finite, there exist $y_i \neq y_j$ with $x_i$ and $x_j$ lying in the same coset of $Z(D^x)$. Therefore, $x + y_i = z(x + y_j)$ for some $z \in Z(D^x)$. As $y_i \neq y_j$, we have $z \neq 1$. Consequently, $(1 - z)x = zy_j - y_i$ lies in $Z(D)$ with $1 - z$ nonzero, hence invertible. Therefore, $x$ is an element in $Z(D)$ and the result is proven. \hfill \Box

Lemma 105.16. Let $D$ be a division ring with center $F$ and $C$ a conjugacy class in $D$ such that $m_F(C)$ is quadratic in $F[t]$. Suppose $f \in D[t]$ has two right roots in $C$. Then $f \in D[t]m_F(C)$ and $f(C) = 0$.

We leave the proof of this as an exercise.

Proposition 105.17. Let $F$ be a real closed field and $H$ the Hamilton quaternions over $F$. Let $f \in H[t]$ be nonzero. Then the following are equivalent:

(1) The polynomial $f$ has infinitely many roots in $D$.
(2) There exist elements $a$ and $b$ in $F$ with $b$ nonzero satisfying $f(a + bi) = 0 = f(a - bi)$. 
(3) The polynomial \( f \) has a right factor \( q(t) \) in \( D[t] \) with \( q \in F[t] \) a quadratic irreducible.

If these conditions hold, then \( f \) vanishes on the conjugacy class of \( a + bi \).

**Proof.** (1) \( \Rightarrow \) (3): By Proposition 105.11, \( f \) has two roots in some conjugacy class \( C \) in \( \mathcal{H} \). Then \( m_F(C) \) lies in \( F[t] \) and is an irreducible quadratic. By Lemma 105.16, we have \( f \) lies in \( \mathcal{H}[t]m_F(C) \) (and \( m_F(C) \) vanishes on \( C \)).

(3) \( \Rightarrow \) (1): Let \( c \) be a root of \( m_F(C) \) in \( F(i) \). Then by Theorem 105.15, \( x \) has infinitely many conjugates in \( \mathcal{H} \). Each of these is a root of \( m_F(C) \) hence of \( f \).

(3) \( \Rightarrow \) (2) is clear and (2) \( \Rightarrow \) (3) follows from Lemma 105.1 and Proposition 105.2.

**Exercises 105.18.**

1. Let \( D \) be a division ring with center \( F \) and \( a, b \) in \( D \) both algebraic over \( F \). Then \( a \) and \( b \) are conjugate in \( D \) if and only if \( m_F(a) = m_F(b) \).

2. Prove Lemma 105.16.

3. Let \( R \) be a real closed field and \( \mathcal{H} \) the Hamilton quaternions over \( F \). Suppose that \( f = \sum_{i=0}^{n} a_i t^i \in \mathcal{H}[t] \) with \( a_0 \in \mathcal{H} \setminus F \) and \( a_1, \ldots, a_n \in F \). Show that \( f \) has at most \( n \) roots in \( \mathcal{H} \). In particular, show \( t^n - a_0 \) has exactly \( n \) solutions in \( D \), and they all lie in \( F(a_0) \).
CHAPTER XIX

Introduction to Representation Theory

In this chapter, we give applications of the Artin-Wedderburn Theorem to finite group theory. In particular, we study group homomorphisms \( \varphi : G \to \text{GL}_n(F) \) with \( G \) a finite group and \( F \) a field, especially the case that the characteristic of \( F \) is zero. We give two important applications, the first a famous theorem of Hurwitz on the products of sums of squares, and the second the theorem of Burnside showing that groups of order \( p^aq^b \) are solvable where \( p \) and \( q \) are primes and \( a, b \) non-negative integers. To prove the second we introduce and study characters of representations, i.e., the trace of a representation \( \varphi : G \to \text{GL}_n(F) \).

Throughout this chapter \( F \) will denote a field and \( G \) a group.

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Throughout this section \( R \) will denote a nonzero commutative ring.

**Definition 106.1.** Let \( R \) be a commutative ring and \( G \) a group (respectively, monoid) (usually written multiplicatively). The group ring (respectively, monoid ring) \( R[G] \) of \( G \) over \( R \) is the free \( R \)-module on basis \( G = \{ g \mid g \in G \} \) made into a ring by the following multiplication:

\[
(\sum_G a_g g)(\sum_G b_h h) = \sum_G c_k k \quad \text{with} \quad a_h, b_h \in G \quad \text{almost all zero where}
\]

\[
c_k = \sum_{gh = k} a_g b_h.
\]

We have \( 1_{R[G]} = 1_R e_G \) where \( e_G \) is the unity of \( G \).

The map \( \varphi_{R[G]} : R \to R[G] \) given by \( r \mapsto r1_{R[G]} \) is a ring monomorphism, so we view \( R \subset R[G] \). As \( R \subset Z(R[G]) \), the center of \( R[G] \), the ring \( R[G] \) is an \( R \)-algebra.

**Examples 106.2.** Let \( R \) be a commutative ring and \( G \) a group.

1. If \( H = \{ t^i \mid i \geq 0 \} \), then \( R[H] = R[t] \).
2. If \( H = \mathbb{N}^n \) with \( \mathbb{N} = \mathbb{Z}^+ \cup \{0\} \), (an additive monoid), then \( R[H] \cong R[t_1, \ldots, t_n] \).
3. $R[Z] \cong R[t, t^{-1}]$.
4. $R[Z^n] \cong R[t_1, t_1^{-1}, \ldots, t_n, t_n^{-1}]$.
5. $R[G]$ is commutative if and only if $G$ is abelian.
7. If $G$ is cyclic of order $n$, then $R[G] \cong R[t]/(t^n - 1)$.

Remarks 106.3. Let $R$ be a commutative ring, $G$ a (multiplicative) group, $M$ an $R[G]$-module, and $N$ an $R$-module.

1. The map $\lambda : G \to \text{Aut}_R(M)$ given by $g \mapsto \lambda_g = \lambda(g) : m \mapsto gm$, is a group homomorphism. (The inverse of $\lambda_g$ is $\lambda_{g^{-1}}$.)
2. A group homomorphism $\varphi : G \to \text{Aut}_R N$ is called a representation of $G$. Such a $\varphi$ makes $N$ into an $R[G]$-module via
   
   $g \cdot x := \varphi(g)(x)$ for all $g \in G$, $x \in N$.

Conclusion. A representation $\varphi : G \to \text{Aut}_R N$ is equivalent to an $R[G]$-structure on $N$. We say that $\varphi$ is irreducible if $N$ is an irreducible $R[G]$-module (via $\varphi$).

3. If $\tau : G' \to G$ is a group homomorphism, then any representation of $G$ induces a representation of $G'$ by composition.

4. As $G$ is an $R$-basis for the $R$-free module $R[G]$, any representation $\varphi : G \to \text{Aut}_R N$ can be extended to an $R$-algebra homomorphism which we also write as $\varphi : R[G] \to \text{End}_R N$ and also call a representation. Conversely, any $R$-algebra homomorphism $\varphi : R[G] \to \text{End}_R N$ restricts to a representation $\varphi : G \to \text{Aut}_R N$.

Conclusion. A representation $\varphi : G \to \text{Aut}_R N$ is equivalent to an $R$-algebra homomorphism $\varphi : R[G] \to \text{End}_R N$.

5. Let $\lambda : G \to \sum(G)$ be the left regular representation, i.e., $x \mapsto \lambda_x : g \mapsto xg$. Then $\lambda$ induces a representation $\lambda : R[G] \to \text{End}_R(R[G])$ called the (left) regular representation of $G$ (relative to $R$). It is faithful, i.e., $\ker \lambda = 0$.

6. Suppose that $N$ is a finitely generated free $R$-module of rank $n$ on (an ordered) basis $\mathcal{B}$. Then $N \cong R^n$ and we have a group isomorphism $\psi^{-1} : \text{GL}_n(R) \to \text{Aut}_R(N)$ where $\psi(T) = [T]_\mathcal{B}$, the matrix representation of $T$ relative to the basis $\mathcal{B}$. If $\varphi : G \to \text{Aut}_R N$ is a representation, we get a group representation $\tilde{\varphi} = \psi \varphi : G \to \text{GL}_n(R)$ called an $R$-representation of degree $n$. We say that $\tilde{\varphi}$ affords $N$ or $\varphi$ affords $N$ if $\mathcal{B}$ is clear. [It is also common to write $\text{GL}(N)$ for $\text{Aut}_R(N)$.] Hence fixing a basis $\mathcal{B}$ for a finitely generated free $R$-module $N$ gives rise to a representation $\tilde{\varphi} : G \to \text{GL}_n(R)$ where $n$
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is the rank of $N$, hence an $R[G]$-module structure to $R^n$. One often writes $\tilde{\varphi}(n)$, $n \in N$, for $\varphi(g)(n)$ if we know $N$.

7. Suppose that $N$ and $N'$ are two free $R$-modules of the same rank $n$ with $\varphi : G \to \text{Aut}_R(N)$ and $\psi : R \to \text{Aut}_R(N')$ affording $N$ and $N'$, respectively. We say $\varphi$ and $\psi$ are equivalent and write $\varphi \sim \psi$, if there exists an $R$-isomorphism $T : N \to N'$ satisfying

$$
\begin{array}{ccc}
N & T & N' \\
\varphi(g) \downarrow & & \psi(g) \downarrow \\
N & T & N'.
\end{array}
$$

commutes for all $g$ in $G$, i.e., if $N$ is an $R[G]$-module via $\varphi$ and $N'$ is an $R[G]$-module via $\psi$ in the above, then $T : N \to N'$ is an $R[G]$-isomorphism. If $N = R^n = N'$, of course, we write $GL_n(R)$ for $\text{Aut}_R N$.

**Conclusion.** We have $\varphi \sim \psi$ in the above if and only if $N \cong N'$ as $R[G]$-modules.

8. We say the $R[G]$-module $M$ is $G$-trivial if $gm = m$ for all $m$ in $M$ and all $g$ in $G$.

9. Define

$$M^G := \{ m \in M \mid gm = m \text{ for all } g \in G \},$$

a submodule of $M$ called the set of $G$-fixed points of $M$.

10. Let $\varphi : G \to GL_n(R)$ afford the finitely generated free $R$-module $N$. Then $N$ is $G$-trivial via $\varphi$ if and only if $\varphi$ is the trivial map, i.e., $g \mapsto I$ for all $g$ in $G$. We call this the trivial representation of $G$. If $n = 1$, it is irreducible.

11. Suppose both $M$, $N$, and $P$ are $R[G]$-modules. Then $\text{Hom}_R(M, N)$ is an $R[G]$-module defined by

$$(\sigma f)(m) := \sigma(f(\sigma^{-1}(m)))$$

for all $\sigma \in G$, $m \in M$, and $f \in \text{Hom}_R(M, N)$.

**Check.** If $f \in \text{Hom}_R(M, N)$, $g \in \text{Hom}_R(N, P)$, then $\sigma(g \circ f) = (\sigma g) \circ (\sigma f)$ and

$$\text{Hom}_{R[G]}(M, N) = (\text{Hom}_R(M, N))^G.$$

12. We view $R$ as a trivial $R[G]$-module. If $M$ is an $R[G]$-module, then the isomorphism of $R[G]$-modules $\text{Hom}_R(R, M) \to M$ given by $f \mapsto f(1)$ induces an isomorphism $\text{Hom}_{R[G]}(R, M) \cong M^G$. 

13. Suppose that $G$ is a finite group. Define the norm of $G$ by

$$N_G := \sum_{g \in G} g \in R[G].$$

As $\sigma N_G = N_G = N_G \sigma$ for all $\sigma$ in $G$, we have $N_G$ lies in $R[G]^G$ and $N_G x$ lies in $M^G$ for all $x$ in $M$, i.e., $N_G : M \to M^G$ given by $x \mapsto N_G x$ is an $R[G]$-homomorphism, so $N_G M \subset M^G$. The quotient $M^G/N_G M$ is an object of study in algebraic number theory. Of course, this also says that $N_G$ lies in the center of $R[G]$. More generally, if $H$ is a normal subgroup of $G$, then $N_H$ lies in the center of $R[G]$. 

14. Suppose that $G$ is a finite group and both $M$ and $N$ are $R[G]$-modules. Let $f : M \to N$ be an $R$-homomorphism. Define the trace of $f$ to be the $R[G]$-homomorphism

$$\text{Tr}_G f := N_G f = \sum_{G} \sigma f : M \to M.$$

In particular, if $f$ is an $R[G]$-homomorphism, then $\text{Tr}_G f = |G| f$.

**Lemma 106.4.** Let $R$ be a commutative ring and $G$ a finite group. Suppose that $M$, $M'$, $N$, $N'$ are $R[G]$-modules and $f : M \to N$ an $R$-homomorphism. If $\varphi : M' \to M$ and $\psi : N \to N'$ are $R[G]$-homomorphisms, then $\text{Tr}_G (\varphi f \psi) = \psi \circ \text{Tr}_G f \circ \varphi$.

**Proof.** By definition, $(\sigma f)(m) = \sigma(f(\sigma^{-1} m))$, so

$$\text{Tr}_G (\psi f \varphi) = \sum_{G} \sigma \circ (\psi f \varphi) = \sum_{G} (\sigma \circ \psi \circ (\sigma \circ f) \circ (\sigma \circ \varphi)$$

$$= \psi \circ \left(\sum_{G} \sigma f\right) \circ \varphi = \psi \circ \text{Tr}_G f \circ \varphi. \quad \square$$

**Theorem 106.5.** (Maschke’s Theorem) Let $F$ be a field and $G$ a finite group. Suppose that either $\text{char} F = 0$ or $\text{char} F \not| |G|$. Then $F[G]$ is semisimple.

**Proof.** Let $V$ be an $F[G]$-module and $W$ an $F[G]$-submodule. We must show that $W$ is a direct summand of $V$ as an $F[G]$-module. Let $i_W : W \to V$ denote the inclusion map. We must show that this is a split $F[G]$-monomorphism by Exercise 37.19(12). Write $V = W \oplus W'$ as $F$-vector spaces. [Of course, $W'$ need not be an $F[G]$-module.] Let $\pi_W : V \to W$ be the vector space projection of $V$ onto $W$. By assumption, $1/|G| \in F^\times$, so we can define the ‘average’

$$\varphi := \frac{1}{|G|} \text{Tr}_G \pi_W : V \to W.$$
The map \( \varphi \) is an \( F[G] \)-homomorphism and, by the lemma, we have
\[
\text{Tr}_G(i_W \pi_W i_W) = i_W \circ \text{Tr}_G \pi_W \circ i_W = |G| i_W \circ \varphi \circ i_W.
\]
Since
\[
\text{Tr}_G(i_W \pi_W i_W) = \text{Tr}_G (i_W (\pi_W i_W)) = \text{Tr}_G(i_W 1_W) = \text{Tr}_G i_W = |G| i_W,
\]
we conclude that \( i_W \circ (\varphi \circ i_W) = i_W = i_W 1_W \) as \( |G| \) is a unit in \( F[G] \).
As \( i_W \) is a monomorphism, we have \( \varphi \circ i_W = 1_W \), i.e., \( i_W \) is a split \( F[G] \)-monomorphism as desired. \( \square \)

Maschke’s Theorem allows us to apply the Artin-Wedderburn Theorem to \( F[G] \) when \( G \) is a finite group and \( \text{char } F = 0 \) or \( \text{char } F \) does not divide \( |G| \). We also get one further piece of useful information, as we can compute the dimension of the center of \( F[G] \), which we now do.

**Definition 106.6.** Let \( R \) be a commutative ring and \( G \) a group. If \( g \) is an element of \( G \) and the conjugacy class \( C(g) \) of \( g \) in \( G \) finite, we let
\[
C_g := \sum_{C(g)} h \text{ called the class sum of } g \text{ in } R[G].
\]

**Lemma 106.7.** Let \( R \) be a commutative ring and \( G \) a finite group. If \( C_{g_1}, \ldots, C_{g_r} \) are the distinct class sums of \( G \) in \( R[G] \), then the center \( Z(R[G]) \) is a free \( R \)-module on basis \( \mathcal{B} := \{C_{g_1}, \ldots, C_{g_r}\} \). In particular, the rank of \( Z(R[G]) \) is equal to the number of conjugacy classes of \( G \).

**Proof.** Let \( \sigma \) be an element of \( G \). Then for all \( g \) in \( G \), we have \( \sigma C_g \sigma^{-1} = C_g \), so \( C_g \) lies in \( Z(R[G]) \). Since \( G = \bigvee C(g_i) \) is an \( R \)-basis for \( R[G] \), the set \( \mathcal{B} \) is \( R \)-linearly independent. Thus we need only show that \( \mathcal{B} \) spans \( Z(R[G]) \). Let \( y = \sum a_g g \) lie in \( Z(R[G]) \). Then for each \( \sigma \) in \( G \), we have \( y = \sigma y \sigma^{-1} = \sum C_g \). As \( G \) is an \( R \)-basis for \( R[G] \), we have \( a_g = a_{g \sigma} \) for all \( \sigma \) in \( G \), and the result follows. \( \square \)

We now apply the Artin-Wedderburn Theorem to the case that \( F \) is an algebraically closed field and \( G \) is a finite group with \( \text{char } F = 0 \) or \( \text{char } F \nmid |G| \), i.e., Maschke’s Theorem holds, e.g., \( F = \mathbb{C} \), to conclude the following:

**Theorem 106.8.** Let \( G \) be a finite group and \( F \) an algebraically closed field of characteristic zero or positive characteristic not dividing the order of \( G \). Let \( \mathfrak{A}_1, \ldots, \mathfrak{A}_r \) be a basic set for \( F[G] \) and \( n_i = \dim_F \mathfrak{A}_i \), \( i = 1, \ldots, r \). Then
1. \( F[G] \cong \bigsqcup_{i=1}^r \mathfrak{A}_i^{n_i} \) as \( F[G] \) modules.
2. \( F[G] \cong \bigotimes_{i=1}^r \mathfrak{A}_i^{n_i}(F) \) as rings.
3. \( r \) is the number of conjugacy classes of \( G \).
4. \( |G| = \sum_{i=1}^r n_i^2 \) and at least one of the \( n_i \) is one.
Proof. There are no finite dimensional $F$-division rings over $F$ except for $F$, as any such would contain an element $x$ not in $F$, but then $F(x)$ would be a commutative division ring distinct from $F$. It also follows that if $B_i$ is the Wedderburn component corresponding to $\mathbb{M}_{n_i}(F)$, we have $\text{dim}_F B_i = n_i$ is the dimension of a column space of $\mathbb{M}_{n_i}(F)$. Therefore, we have (1), (2), and (3). Finally, $Z(M_{n_i}(F)) = F$, so

$$G = \text{dim}_C C[G] = \text{dim}_C \left( \times_{i=1}^r \mathbb{M}_{n_i}(C) \right) = \sum_{i=1}^r n_i^2,$$

and the trivial representation of $G$ into $\text{GL}_1(F) = F^\times$ is irreducible of degree one giving (4). $\square$

We can weaken the hypothesis of the theorem that $F$ be algebraically close to having all the Wedderburn components matrix rings over $F$. We shall investigate this in the next section.

Remarks 106.9. Let $G$ be a finite group and $F$ a field of characteristic zero or of positive characteristic not dividing $|G|$.

1. A representation $\varphi : G \to \text{GL}_1(F)$, i.e., of degree one is called a linear representation). By Maschke’s Theorem, it must be irreducible.

2. As a special case of (1), suppose that $G$ is an elementary 2-group, i.e., $G \cong (\mathbb{Z}/2\mathbb{Z})^n$ for some $n$ and $F$ is a field of characteristic different from two. Then Maschke’s Theorem holds. There exist at least $2^n$ linear inequivalent representations of $G$ given by the trivial representation and the $2^n - 1$ representations each defined by taking one nonzero element of $G$ to $-1$ and all other elements if $G$ to $1$. Since $\text{dim}_F F[G] = 2^n$ and $F[G]$ is semisimple, these must be all the irreducible representations of $G$.

3. Let $G^{ab} = G/[G,G]$. Then any representation of $G^{ab}$ gives rise to a representation of $G$ of the same degree via composition with the canonical epimorphism $\pi : G \to G^{ab}$.

4. If $F$ is algebraically closed, then any irreducible representation $\varphi : G^{ab} \to \mathbb{M}_n(F)$ must be of degree one.

Example 106.10. Let $G$ be a finite group of order $n$. We give some examples of complex representations, i.e., representations $\varphi : G \to \text{GL}_n(C)$.

1. We interpret Theorem 106.8 in the case of complex representations. Let $\mathfrak{A}_1, \ldots, \mathfrak{A}_r$ be a basic set for $C[G]$ and $n_i = \text{dim}_C \mathfrak{A}_i$, $i = 1, \ldots, r$. Let $\varphi_i$ be the irreducible representation afforded by $\mathfrak{A}_i$. Then the regular representation $\lambda : G \to \text{GL}_n(C)$ is similar to
2. Let \( G = \langle x \rangle \) be a cyclic group of order \( n \) and \( \omega \) a primitive root of unity in \( \mathbb{C} \). Then \( \rho_i : G \to \mathbb{C}^\times = \text{GL}_1(\mathbb{C}) \) determined by \( x \to \omega^i \) for \( i = 1, \ldots, n \) are all the inequivalent irreducible representations of \( G \).

3. Let \( G \) be the dihedral group \( D_3 \) with \( G = \{a, b \mid a^3 = 1 = b^2, bab^{-1} = a^{-1}\} \). We know the trivial representation is of degree one, and there are three conjugacy classes in \( D_3 \). Therefore, there is another a linear character and one of degree 2. The non-trivial linear character \( G^{ab} \cong \mathbb{Z}/2\mathbb{Z} \) induces the irreducible character \( G \to \mathbb{C}^\times \) determined by \( a \mapsto 1 \) and \( b \mapsto -1 \). Let \( \omega \) be a primitive cube root of unity. The representation \( \varphi : G \to \text{GL}_n(\mathbb{C}) \) determined by

\[
\begin{pmatrix}
\varphi_3 & n_1 \\
\vdots \\
\varphi_3 & n_r
\end{pmatrix}
\]

is the third irreducible representation.

**Exercises 106.11.**

1. Let \( F \) be an algebraically closed field of characteristic zero, \( G \) be a finite group and \( \varphi : G \to \text{GL}_n(F) \) a representation. Show that \( \varphi \) is irreducible if and only if for each \( x \in G \), there exists an element \( \lambda \in F \) such that \( \lambda(x) = \lambda I \).

2. Let \( F \) be an algebraically closed field of characteristic zero, \( G \) be a finite group, and \( z \) an element of the center of \( F[G] \). Show if \( V \) is an irreducible \( F[G]\)-module, then there exists an element \( \lambda \in F \) satisfying \( zv = \lambda v \) for all \( v \in V \).

3. Let \( F \) be an algebraically closed field of characteristic zero and \( G \) be a finite group. Show if there exists a faithful irreducible \( F[G]\)-module \( M \) (i.e., irreducible and \( gm = m \) for all \( m \in M \), then \( g = e_G \)), then the center of \( G \) is cyclic.

4. Determine all complex representations of a finite abelian group.
5. Let $G = \langle x \rangle$ be a cyclic group of order $n$ and $\omega$ a primitive root of unity in $\mathbb{C}$. In the notation of Example 106.10(2) determine the $\mathbb{C}[G]$-submodule of $\mathbb{C}[G]$ determined by $\rho_1$, by $\rho_i$.

6. Let $G$ be the dihedral group $D_3$ and $\omega$ a primitive root of unity in $\mathbb{C}$. In the notation of Example 106.10(2) determine the irreducible $\mathbb{C}[G]$-submodules of $\mathbb{C}[G]$.

7. Find all complex irreducible representations of the dihedral group $D_4$.

8. Find all complex irreducible representations of the quaternion group $Q_3$.

9. Find all real irreducible representations of the quaternion group $Q_3$.

107. Split Group Rings

The purpose of this section is to extend theorems about of group rings over fields with the group finite and Maschke’s Theorem holds when the underlying field is algebraically closed to the case that the underlying division rings in the Wedderburn decomposition of a group ring are always the base field. This is useful in the general theory.

**Definition 107.1.** Let $D$ be a division ring and $M$ a (left) $D$-vector space. If $A$ is a subring of $\text{End}_D(M)$, we say that $A$ act densely on $M$ if given any $D$-linear independent vectors $v_1, \ldots, v_m$ in $M$, $m > 0$, and vectors $v'_1, \ldots, v'_n$ in $M$, then there exists a $T \in A$ satisfying $Tv_i = v'_i$ for $i = 1, \ldots, n$.

**Theorem 107.2.** (Jacobson Density Theorem) Let $R$ be a nonzero commutative ring, $M$ an irreducible $R$-module, and $D = \text{End}_R(M)$ (a division ring by Schur’s Lemma). Then

$$\lambda : R \to \text{End}_D(M) \text{ by } a \mapsto \lambda_a : m \mapsto am$$

is a ring homomorphism. Set

$$L_R(M) := \text{im } \lambda \subset \text{End}_D(M).$$

Then $L_R(M)$ acts densely on $M$ as a $D$-vector space. In particular, if $\text{dim}_D M$ is finite, then $\lambda$ is a ring epimorphism.

**Proof.** Let $a$ be an element on $R$ and $f$ an element of $\text{End}_D(M)$. Then for all $m$ in $M$, we have

$$\lambda_a f(m) = af(m) = f(am) = f\lambda_a(m),$$

so $\lambda_a$ lies in $\text{End}_D(M)$ for all $a$ in $R$, and clearly, $\lambda$ is a ring homomorphism. Let $v_1, \ldots, v_m$ in $M$ be $D$-linearly independent with $m \geq 1$ and $v'_1, \ldots, v'_n$ in $M$. As $M$ is $R$-irreducible, $M = Rv_i$ for every $i = 1, \ldots, n$. 
Thus the result is trivial if $m = 1$. So assume that $m > 1$. We are finished if we can establish the following:

**Claim:** It suffices to produce $s_m \in R$ satisfying

$$s_m v_i = \begin{cases} 0, & \text{if } i < m, \\ \neq 0, & \text{if } i = m. \end{cases}$$

Indeed, suppose that such an $s_m$ exists. As $s_m v_m \neq 0$, we have $M = Rs_m v_m$ by irreducibility. Write $v'_m = rs_m v_m$ with $r \in R$ and let $b_m = rs_m$. In an analogous way, there exist $b_j \in R$ satisfying $b_j v_j = \delta_{ij} v'_j$, with $\delta_{ij}$ the Kronecker delta. Then $b = b_1 + \cdots + b_m$ works. This establishes the claim.

Now assume that the result is false. This means, we have

(*) $\text{If } a \in R \text{ satisfies } av_i = 0, 1 \leq i \leq m - 1, \text{ then } av_m = 0.$

By induction, $M^{m-1} = \{ (av_1, \ldots, av_{m-1}) \mid a \in R \}$, so (*) implies if $a, b \in R$, then

$$(av_1, \ldots, av_{m-1}) = (bv_1, \ldots, bv_{m-1}) \Rightarrow av_m = bv_m.$$ 

This means that

$$\mu : M^{m-1} \to M \text{ given by } (av_1, \ldots, av_{m-1}) \mapsto av_m$$

is a well-defined $R$-homomorphism. Let $\iota_j : M \to M^{m-1}$ be the $R$-homomorphism given by $m \mapsto (0, \ldots, m \underset{j}{\ldots}, 0, \ldots, 0)$ as usual, and set

$$\varphi_j = \mu \iota_j \text{ in } D = \text{End}_R(M).$$

Then

$$v_m = \mu(av_1, \ldots, av_{m-1}) = \sum_{j=1}^{m-1} \mu(0, \ldots, m \underset{j}{\ldots}, 0, \ldots, 0) = \sum_{j=1}^{m-1} \varphi_j(v_j),$$

contradicting $v_1, \ldots, v_m$ are $D$-linearly independent. \hfill \Box

**Remark 107.3.** Let $F$ be a field and $A$ a finite dimensional $F$-algebra (i.e., $\dim_F A < \infty$). Let $M$ be a finitely generated $A$-module, hence a finite dimensional $F$-vector space. Then $\text{End}_A(M) \subset \text{End}_F(M)$ is a subring. By definition $F \subset Z(A)$, the center of $A$, so

$$\lambda : A \to \text{End}_F(M) \text{ by } a \mapsto \lambda_a : m \mapsto am$$

is a ring homomorphism. Set

$$L_A(M) := \text{im } \lambda \cong A/\ker \lambda = A/\text{ann}_R M.$$ 

We view

$$F \subset \text{End}_A(M) \subset \text{End}_F(M)$$

as subrings via $a \mapsto a1_M$. 

Theorem 107.4. (Burnside) Let \( F \) be a field, \( A \) an \( F \)-algebra (not necessarily finitely generated), and \( M \) an irreducible \( A \)-module that is also a finite-dimensional vector space over \( F \) (so cyclic as finitely generated).

(1) If \( F \) is algebraically closed, then \( F = \text{End}_A(M) \).
(2) If \( F = \text{End}_A(M) \), then \( L_A(M) = \text{End}_F(M) \).

Proof. \( \text{End}_A(M) \) is a finite dimensional vector space over \( F \) as \( M \) is. By Schur’s Lemma 99.15, \( D = \text{End}_A(M) \) is a division ring.

(1): Let \( x \in D \). Then \( F \subset Z(D) \), so \( F(x) \subset D \) is a commutative division ring, i.e., \( F(x) \) is a field. As \( [F(x) : F] \leq \dim_F D < \infty \), the field extension \( F(x)/F \) is algebraic, hence \( F(x) = F \), as \( F \) is algebraically closed.

(2): As \( \dim_R(M) \) is finite, this follows from the Jacobson Density Theorem 107.2. \( \square \)

Definition 107.5. Let \( A \) be a semi-simple ring and a finite dimensional \( F \)-algebra with \( F \) a field. Suppose that \( A = B_1 \oplus \cdots \oplus B_m \) is a Wedderburn decomposition. We say that \( A \) is \( F \)-split if \( B_i \) is a matrix ring over \( F \) for every \( i = 1, \ldots, m \).

We leave the following as an exercise:

Proposition 107.6. Let \( F \) be a field, \( A \) a semi-simple finite dimensional \( F \)-algebra. Then \( A \) is \( F \)-split if and only if \( F \cong \text{End}_A(M) \) for every irreducible \( A \)-module \( M \). In particularly, if \( F \) is algebraically closed, then \( A \) is \( F \)-split.

Remark 107.7. Suppose that \( F \) is a field, \( G \) a finite group not divisible by the characteristic of \( F \). Then \( F[G] \) is semi-simple. Suppose, in addition, that \( F[G] \) is \( F \)-split and \( M_1, \ldots, M_r \) a complete set of representatives for the isomorphism classes of irreducible \( F[G] \)-modules. Then \( Z(F[G]) = \prod_{i=1}^r F \) and \( r \) is the number of conjugacy classes of \( G \) by Lemma 106.7.

Using the Artin-Wedderburn Theorem, Maschke’s Theorem, Burnside’s Lemma, and the above, we now have the following generalization of Theorem 106.8:

Summary 107.8. Let \( F \) be a field, \( G \) a finite group such that \( \text{char } F \nmid |G| \) (so \( F[G] \) is semi-simple by Maschke’s Theorem). Suppose that \( F[G] \) is \( F \)-split, e.g., if \( F \) is algebraically closed. Let \( M_1, \ldots, M_r \) be a complete set of representatives for the isomorphism classes of irreducible \( F[G] \)-modules and \( n_i = \dim_F M_i \) for \( 1 \leq i \leq r \). Then
108. ADDENDUM: HURWITZ’S THEOREM 751

1. $F[G] \cong \prod_{i=1}^{r} M_{n_i}^i$ as $F[G]$-modules.
2. $F[G] \cong \times_{i=1}^{r} M_{n_i}(F)$ as rings.
3. $r$ is the number of conjugacy classes of $G$.
4. $|G| = \sum_{i=1}^{r} n_i^2$.
5. $L_{F[G]}(M_i) = \text{End}_F(M_i)$ for $i = 1, \ldots, r$.

Exercises 107.9.
1. Let $F$ be a field, $A$ an $F$-algebra (not necessarily finitely generated), and $M$ a completely reducible $A$-module. Show if $\text{End}_A(M) \cong F$, then $M$ is an irreducible $A$-module.

108. Addendum: Hurwitz’s Theorem

We shall give a nice application of the theory developed so far. If $n \geq 3$, the general quaternion group $Q_n$ is the group on $n$ generators $a_1, \ldots, a_{n-1}, \varepsilon$ subject to the relations

\[
\varepsilon^2 = 1 \quad a_i^2 = \varepsilon \quad \text{for each } i = 1, \ldots, n-1 \quad a_i a_j = \varepsilon a_j a_i \quad \text{for all } i, j = 1, \ldots, n-1 \text{ satisfying } i \neq j.
\]

For example, $Q_3$ is the usual quaternion group. We shall see below that this group has order $2^n$.

Properties 108.1. Let $Q_n$, $n \geq 3$, be the general quaternion group.
1. The element $\varepsilon$ lies in the center of $Q_n$.
2. The element $\varepsilon$ satisfies $\varepsilon = a_i a_j a_i^{-1} a_j^{-1}$ for all $i \neq j$. In particular, $\varepsilon$ lies in $[Q_n, Q_n]$.
3. $[Q_n, Q_n] = \langle \varepsilon \rangle = \{1, \varepsilon\}$ (as $Q_n/\langle \varepsilon \rangle$ is abelian and $\varepsilon \in [Q_n, Q_n]$). In particular, $Q_n/\langle \varepsilon \rangle$ is an elementary 2-group on the image of the generators $a_1, \ldots, a_{n-1}$. Hence $|Q_n/\langle \varepsilon \rangle| = 2^{n-1}$ and $|Q_n| = 2^n$.
4. There exist (at least) $2^{n-1}$-linear representations of $Q_n$ over any field of characteristic different from two by Remark 106.9(2).
5. If $n$ is odd, then the center $Z(Q_n)$ of $Q_n$, is $\langle \varepsilon \rangle$:

It suffices to show that $z = a_1 \cdots a_l$ for $1 \leq l \leq n$ is not central. But if such a $z$ lies in $Z(Q_n)$, then

\[
a_1 a_2 \cdots a_l = a_1 a_l a_1 = a_2 a_l a_1 = \ldots = a_l a_1 a_l,
\]

so $l - 1$ is even, hence $l$ is odd and $l < n - 1$. But

\[
a_{n-1} a_1 \cdots a_l = a_1 a_l a_{n-1} = a_l^\varepsilon a_{n-1} a_1 \cdots a_l,
\]

so $l - 1$ is even.
6. If \( n \) is even, then \( Z(\mathbb{Q}_n) = \langle \varepsilon, a_1 \cdots a_{n-1} \rangle \):

This is similar to the previous argument, using \( a_1 \cdots a_n \) is central as \( n-1 \) is odd, so

\[
a_1 \cdots a_{n-1} a_j = \varepsilon^{n-2} a_j a_1 \cdots a_{n-1} = a_j a_1 \cdots a_{n-1},
\]

since \( a_j \) occurs once in \( a_1 \cdots a_{n-1} \).

7. If \( g \) in \( \mathbb{Q}_n \) is not central, then its conjugacy class \( C(g) = \{g, \varepsilon g\} \):

As \( g \notin Z(\mathbb{Q}_n) \), we have \(|C(g)| > 1\). Let \(-: \mathbb{Q}_n \to \mathbb{Q}_n/\langle \varepsilon \rangle\) be the canonical surjection. As \( \mathbb{Q}_n \) is abelian, \( g_o \in \{g, \varepsilon g\} \) if \( g_o = xgx^{-1} \) with \( x \in \mathbb{Q}_n \).

We next look at the degrees of the irreducible representations of the general quaternion group over an algebraically closed field of characteristic different from two.

**Calculation 108.2.** Let \( \mathbb{Q}_n, n \geq 3 \), be the general quaternion group and \( r \) the number of conjugacy classes of \( \mathbb{Q}_n \). Let \( F \) be a algebraically closed field of characteristic different from two. As \( \mathbb{Q}_n \) is a 2-group, we know that \( r \) is the number of irreducible \( F \)-representations of \( \mathbb{Q}_n \) and there exist \( 2^{n-1} \) linear \( F \)-representations by the properties of established above. Therefore, using the properties above, we have

\[
2^n = |G| = \sum_{i=1}^{r} d_i^2 \text{ with } d_i = 1 \text{ for } 1 \leq i \leq 2^{n-1}.
\]

**Case 1.** \( n \) is odd:

We have

\[
r = |Z(\mathbb{Q}_n)| + \frac{|\mathbb{Q}_n| - |Z(\mathbb{Q}_n)|}{2} = 2 + \frac{2^n - 2}{2} = 2^{n-1} + 1.
\]

So there exists one further irreducible \( F \)-representation and it is of degree \( d_r \). Since \( 2^n = 1 \cdot 2^{n-1} + d_r^2 \), we have \( d_r = 2^{\frac{n-1}{2}} \).

**Case 2.** \( n \) is even:

We have

\[
r = |Z(\mathbb{Q}_n)| + \frac{|\mathbb{Q}_n| - |Z(\mathbb{Q}_n)|}{2} = 4 + \frac{2^n - 4}{2} = 2^{n-1} + 2,
\]

so there exist two further irreducible \( F \)-representations of degrees \( d_{r-1} \) and \( d_r \) and these satisfy \( 2^n = 1 \cdot 2^{n-1} + d_{r-1}^2 + d_r^2 \), hence \( 2^{n-1} = d_{r-1}^2 + d_r^2 \). Write \( d_{r-1} = 2^k f_1 \) and \( d_r = 2^l f_2 \) with \( f_1, f_2 \) odd. It follows that \( k = l \), hence

\[
2^{n-2k-1} = f_1^2 + f_2^2 \equiv 2 \mod 4.
\]
If either \( f_1 \) or \( f_2 \) is greater than one, we would have \( 4 \mid f_1^2 + f_2^2 \), which is impossible. Consequently, \( d_{r-1} = d_r = 2^{\frac{n-2}{2}} \).

We shall use the calculation above to solve a very nice problem in field theory. Let \( x_1,\ldots,x_n \) and \( y_1,\ldots,y_n \) be variables over a field \( F \). We wish to know when there exists a formula

\[
(*) \quad z_1^2 + \cdots + z_n^2 = (x_1^2 + \cdots + x_n^2)(y_1^2 + \cdots + y_n^2)
\]

over \( F[x_1,\ldots,x_n,y_1,\ldots,y_n] \) with the \( z_i \) bilinear in the \( x_i \)'s and \( y_j \)'s, i.e., bilinear when plugging in any values in \( F \) for the variables. We know that there are such formulae in certain cases, e.g.,

- If \( n = 1 \). (This is trivial.)
- If \( n = 2 \). (Cf. with the norm from \( \mathbb{C} \) to \( \mathbb{R} \).)
- If \( n = 4 \). (Cf. with the norm form of a quaternion algebra.)
- If \( n = 8 \). (Cf. with the norm form Cayley’s octonion algebra—a non-associate algebra, i.e., satisfies all properties of an algebra except associativity, and in addition all nonzero elements have inverses.)

**Theorem 108.3.** (Hurwitz) Let \( F \) be a field of characteristic not two. Then \((*)\) exists if and only if \( n = 1, 2, 4, \text{ or } 8 \).

**Proof.** (Eckmann) If \( n = 1, 2, 4, 8 \), such a formula exists. We have written it down except for the case of \( n = 8 \), which we leave to the reader to look up. Let \( X = (x_1,\ldots,x_n) \) and suppose that \((*)\) holds. Let \( z_i = \sum_{j=1}^n a_{ij}(X)y_j \) for \( i = 1,\ldots,n \). Each of the \( a_{ij}(X) \) must be linear in the \( x_i \)'s. We have

\[
\sum_{i=1}^n z_i^2 = \sum_{i=1}^n (\sum_{j=1}^n a_{ij}(X)y_j)^2
\]

\[
= \sum_{i,j=1}^n a_{ij}(X)^2 y_j^2 + 2 \sum_{i=1}^n \sum_{1 \leq j < k \leq n} a_{ij}(X)a_{ik}(X)y_jy_k.
\]

Using \((*)\) and comparing coefficients, we see that

\[
\sum_{i=1}^n x_i^2 = \sum_{i=1}^n a_{ij}(X)^2 \quad \text{for } j = 1,\ldots,n.
\]

\[
\sum_{i=1}^n a_{ij}(X)a_{ik}(X) = 0 \quad \text{for all } j \neq k.
\]
Let \( A = (a_{ij}(X)) \) in \( \mathbb{M}_n(F[x_1, \ldots, x_n, y_1, \ldots, y_n]) \). We have

\[ A^t A = \left( \sum_{i=1}^{n} x_i^2 \right) I. \]

Write \( A = A_1 x_1 + \cdots + A_n x_n \) with \( A_i \in \mathbb{M}_n(F) \).

Substitute this into \((\dagger)\) and multiply out. Comparing coefficients yields equations:

(1) \( A_i^t A_j + A_j^t A_i = 0 \) for all \( i, j = 1, \ldots, n \) and \( i \neq j \).

(2) \( A_i^t A_i = I \) for all \( i = 1, \ldots, n \).

In particular, the \( A_i \) are orthogonal matrices. We next normalize these equations by setting

\[ B_i = A_i A_i^t \] for \( 1 \leq i \leq n \), hence \( B_n = A_n A_n^t = I \).

The \( B_i \)'s then satisfy:

(1') \( B_i^t B_j + B_j^t B_i = 0 \) for all \( i, j = 1, \ldots, n - 1 \) and \( i \neq j \).

(2') \( B_i^t B_i = I \) for all \( i = 1, \ldots, n \).

As \( B_n = I \), we have \( B_i^t B_n + B_n^t B_i = 0 \) for \( 1 \leq i < n \), hence

\[ B_i^t = -B_i \] for \( i = 1, \ldots, n - 1 \),

so (1') and (2') become

(1'') \( B_i B_j + B_j B_i = 0 \) for all \( i, j = 1, \ldots, n - 1 \) and \( i \neq j \).

(2'') \( B_i^2 = -I \) for all \( i = 1, \ldots, n \).

Let \( \bar{F} \) be an algebraic closure of \( F \) and \( \rho : Q_n \to \text{GL}_n(\bar{F}) \) be the representation induced by \( a_i \mapsto B_i \) for \( i = 1, \ldots, n - 1 \) and \( \varepsilon \mapsto -I \).

Each linear representation of \( Q_n \) must take \( \varepsilon \mapsto 1 \) as \( a_i a_j = \varepsilon a_j a_i \) for \( i \neq j \) and \( F \) is commutative. Since \( \text{char} F \neq 2 \), \( \rho \) is a direct sum of irreducible representations of degree greater than one using Exercise 100.7(1). In particular,

If \( n \) is odd, then \( 2^{\frac{n-1}{2}} \mid n \), hence \( n = 1 \).

If \( n \) is even, then \( 2^{\frac{n-2}{2}} \mid n \).

Since \( 2^{(n-2)/2} > n \) for \( n \geq 10 \), we have \( n \leq 8 \) and we check:

If \( n = 8 \), then \( 2^{\frac{8-2}{2}} \mid 8 \).

If \( n = 6 \), then \( 2^{\frac{6-2}{2}} \not\mid 6 \).

If \( n = 2, 4 \), then \( 2^{\frac{2-2}{2}} \mid n \). \qed
A more general theorem (that we do not prove) is

**Theorem 108.4.** (Hurwitz-Radon Theorem) Let $x_1, \ldots, x_l$ and $y_1, \ldots, y_n$ be variables over a field $F$ of characteristic not two. Then there exists a formula

$$z_1^2 + \cdots + z_n^2 = (x_1^2 + \cdots + x_l^2)(y_1^2 + \cdots + y_n^2)$$

over $F[x_1, \ldots, x_l, y_1, \ldots, y_n]$ with the $z_i$ bilinear in the $x_i$’s and $y_j$’s if and only if $n = 2^{4\alpha + \beta}$ with $l \leq 8\alpha + 2\beta$.

If one does not insist on the $z_i$ being bilinear in the $x_i$’s and $y_j$’s and works in $F(x_1, \ldots, x_n, y_1, \ldots, y_n)$, then Pfister showed that the product of two sums of $2^n$ squares is a sum of $2^n$ squares.

**Exercise 108.5.** Find all complex irreducible representations of the general quaternion group $Q_n$, $n \geq 3$.

### 109. Characters

Throughout this section $F$ will be a field and $G$ a group.

**Definition 109.1.** Let $\varphi : G \to \text{GL}_n(F)$ be a representation. The map $\chi_\varphi : G \to F = M_1(F)$ defined by

$$\chi_\varphi(x) := \text{trace } \varphi(x) \text{ for all } x \in G$$

is called the **character** of the representation $\varphi$. The degree of the representation $\varphi$ is also called the **degree** of $\chi_\varphi$, so it is the integer $\chi_\varphi(e_G)$ which we shall denote by $\chi_\varphi(1)$, i.e., write $e_G$ as 1. In particular, $\chi_\varphi$ is not a group homomorphism in general. We say that $\chi_\varphi$ is an **irreducible character** if $\varphi$ is an irreducible representation. For example, the character associated to the trivial representation called the **trivial character** is an irreducible character of degree one.

**Remarks 109.2.** Let $V$ be an $F[G]$-module of dimension $n$ as a vector space over $F$ and $\varphi = \varphi_V : G \to \text{GL}_n(F)$ a representation affording $V \cong F^n$ relative to some fixed basis. If $\psi : G \to \text{GL}_n(F)$ affords $V$ relative to another basis, then $\varphi(x)$ and $\psi(x)$ are conjugate for all $x$ in $V$, hence $\chi_\varphi(x) = \chi_\psi(x)$ for all $v$ in $V$. Therefore, $\chi_\varphi$ depends only on the equivalence class of $\varphi$. We often write $\chi_V$ for $\chi_\varphi$. As $F[G]$ is $F$-free on basis $G$, the character $\chi_\varphi$ induces an $F$-linear functional $F[G] \to F$ that we also denote by $\chi_\varphi$.

**Proposition 109.3.** Let $V$, $V'$, and $V''$ be $F[G]$-modules, finite dimensional over $F$. Then

1. $\chi_V(\sigma x \sigma^{-1}) = \chi_V(x)$ for all $\sigma$ in $G$. 


If $0 \to V' \to V \to V'' \to 0$ is an exact sequence of finite dimensional $F[G]$-modules, then

$$\chi_V = \chi_{V'} + \chi_{V''}.$$ 

**Proof.** We already showed (1). As for (2), we may assume that $V' \subset V$ with $B'$ a basis for $V'$ extended to a basis $B$ for $V$. 

If $\varphi_V : G \to \text{GL}_n(F)$, then with the obvious notation, 

$$\varphi_V(x) = \begin{pmatrix} \varphi_{V'}(x) & * \\ 0 & \varphi_{V''}(x) \end{pmatrix},$$

and the result follows. $\Box$

**Example 109.4.** Let $G$ be a finite group with char $F$ not dividing $|G|$, so $F[G]$ is semi-simple by Maschke’s Theorem 106.5. Let 

$$\lambda : G \to \text{Aut}_F(F[G])$$ 

be given by $x \mapsto \lambda_x : y \mapsto xy$. 

Extend this representation linearly to an $F$-algebra homomorphism 

$$\lambda : F[G] \to \text{End}_F(F[G]).$$ 

This is just the (left) regular representation. 

Let \{A_1, \ldots, A_r\} be the basic set for $F[G]$. 

Suppose that $F[G]$ is $F$-split, e.g., if $F$ is algebraically closed, and $\varphi_{A_1}, \ldots, \varphi_{A_r}$ are all the irreducible representations of $G$ afforded by the $A_i$ with the $\chi_{A_i}$ the corresponding characters. Set $n_i = \deg \chi_{A_i} = \chi_{A_i}(1)$. Then we have 

$$\chi_{F[G]} = \chi_\lambda = \sum_{i=1}^r n_i \chi_{A_i} = \sum_{i=1}^r \chi_{A_i}(1) \chi_{A_i}.$$ 

Note: $\lambda_x$ with $x \neq e_G$ in $G$ permutes $G$ without fixed points and $G$ is a basis for $F[G]$, so 

$$\text{trace } \varphi_\lambda(x) = \chi_\lambda(x) = 0$$ 

for all $x \neq e_G$. 

More generally, if $\varphi : G \to \text{GL}_n(F)$ is a representation affording a finite dimensional $F$-vector space $V$, then there exist unique integers $m_i \geq 0$ satisfying 

1. $V \cong \prod A_i^{m_i}$ as $F[G]$-modules.
2. $\chi_V = \sum m_i \chi_{A_i}$.
3. $\chi_V(1) = \dim_F V$.

We now show that if $F$ is a field of characteristic zero, then the equivalence class of a representation of a finite group $G \to \text{GL}_n(F)$ is completely determined by its character. This is false if $F$ has positive characteristic dividing the order of $G$. 

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**Proof.** We already showed (1). As for (2), we may assume that $V' \subset V$ with $B'$ a basis for $V'$ extended to a basis $B$ for $V$. 

If $\varphi_V : G \to \text{GL}_n(F)$, then with the obvious notation, 

$$\varphi_V(x) = \begin{pmatrix} \varphi_{V'}(x) & * \\ 0 & \varphi_{V''}(x) \end{pmatrix},$$

and the result follows. $\Box$ 

**Example 109.4.** Let $G$ be a finite group with char $F$ not dividing $|G|$, so $F[G]$ is semi-simple by Maschke’s Theorem 106.5. Let 

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be given by $x \mapsto \lambda_x : y \mapsto xy$. 

Extend this representation linearly to an $F$-algebra homomorphism 

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This is just the (left) regular representation. 

Let \{A_1, \ldots, A_r\} be the basic set for $F[G]$. 

Suppose that $F[G]$ is $F$-split, e.g., if $F$ is algebraically closed, and $\varphi_{A_1}, \ldots, \varphi_{A_r}$ are all the irreducible representations of $G$ afforded by the $A_i$ with the $\chi_{A_i}$ the corresponding characters. Set $n_i = \deg \chi_{A_i} = \chi_{A_i}(1)$. Then we have 

$$\chi_{F[G]} = \chi_\lambda = \sum_{i=1}^r n_i \chi_{A_i} = \sum_{i=1}^r \chi_{A_i}(1) \chi_{A_i}.$$ 

Note: $\lambda_x$ with $x \neq e_G$ in $G$ permutes $G$ without fixed points and $G$ is a basis for $F[G]$, so 

$$\text{trace } \varphi_\lambda(x) = \chi_\lambda(x) = 0$$ 

for all $x \neq e_G$. 

More generally, if $\varphi : G \to \text{GL}_n(F)$ is a representation affording a finite dimensional $F$-vector space $V$, then there exist unique integers $m_i \geq 0$ satisfying 

1. $V \cong \prod A_i^{m_i}$ as $F[G]$-modules.
2. $\chi_V = \sum m_i \chi_{A_i}$.
3. $\chi_V(1) = \dim_F V$.

We now show that if $F$ is a field of characteristic zero, then the equivalence class of a representation of a finite group $G \to \text{GL}_n(F)$ is completely determined by its character. This is false if $F$ has positive characteristic dividing the order of $G$. 

---

**Proof.** We already showed (1). As for (2), we may assume that $V' \subset V$ with $B'$ a basis for $V'$ extended to a basis $B$ for $V$. 

If $\varphi_V : G \to \text{GL}_n(F)$, then with the obvious notation, 

$$\varphi_V(x) = \begin{pmatrix} \varphi_{V'}(x) & * \\ 0 & \varphi_{V''}(x) \end{pmatrix},$$

and the result follows. $\Box$
Theorem 109.5. Let $F$ be a field of characteristic zero, $G$ a finite group, and $V, V'$ two finitely generated $F[G]$-modules. Then

$$V \cong V' \text{ if and only if } \chi_V = \chi_{V'}.$$  

Proof. ($\Rightarrow$) is trivial.

($\Leftarrow$): We know that the $F$-algebra homomorphism $\lambda : F[G] \to \text{End}_F(V)$ induced by $x \mapsto \lambda_x : v \mapsto xv$ affords $V$ and that $F[G]$ is semi-simple by Maschke’s Theorem. Let $\{\mathfrak{A}_1, \ldots, \mathfrak{A}_r\}$ be a basic set and $B_i = B_{\mathfrak{A}_i}$, $1 \leq i \leq r$ the simple components, so $F[G] = B_1 \oplus \cdots \oplus B_r$ is the Wedderburn decomposition. Let $e_i = 1_{B_i}$, $1 \leq i \leq r$, so $e_1, \ldots, e_r$ are central orthogonal idempotents. $V$ is a finitely generated $F[G]$-module, so there exist unique $m_i \geq 0$ satisfying $V \cong \bigoplus A_i^{m_i}$ as $F[G]$-modules, which we view as an equality. 

To prove the theorem, we need to show that the $F$-linear functional $\chi_V : F[G] \to F$ determines all the $m_i$. Fix a $j$, $1 \leq j \leq r$, and let $d_j = \dim_F \mathfrak{A}_j \neq 0$. [Note: We are not assuming $F[G]$ is $F$-split, so the value $d_j$ is not so clear.] As $e_j$ is central, for each $i$, the map $\lambda_{e_j} : \mathfrak{A}_i \to \mathfrak{A}_j$ given by $a \mapsto e_j a$ is an $F[G]$-homomorphism. Viewing $\mathfrak{A}_i$ as a $B_i$-module, we have $\lambda_{e_j} = \delta_{ij} 1_{B_i}$. Thus we can view $\lambda_{e_j} : V \to V$ by $v \mapsto e_j v$ as an $F[G]$-homomorphism and conclude that $\text{im} \lambda_{e_j} = \lambda_{e_j}(V) = \mathfrak{A}_j^{m_j}$. Let $B_{k_j}$ be an (ordered) $F$-basis for the $k$th component of $\mathfrak{A}_j^{m_j}$. Then $B_j = B_{k_1} \cup \cdots \cup B_{k_{m_j}}$ is an $F$-basis for $\mathfrak{A}_j^{m_j}$ and $\mathcal{B} = \bigcup B_j$ is an $F$-basis for $V$. We then have $[\lambda_{e_j}]_{\mathcal{B}}$ is the matrix

$$
\begin{pmatrix}
0 \\
\ddots \\
1 & \cdots & m_j \\
\ddots \\
1 \\
\cdots \\
0
\end{pmatrix}
$$

with $I$ the $d_j \times d_j$ identity matrix. Consequently,

$$\chi_V(e_j) = \text{trace}[\lambda_{e_j}]_{\mathcal{B}} = d_j m_j,$$

so

$$m_j = \frac{\chi_V(e_j)}{d_j} = \frac{\chi_V(e_j)}{\dim_F \mathfrak{A}_j} \text{ in } F$$

is determined by $\chi_V$. \qed
Exercise 109.6. Let $G$ be a finite $p$-group and $F$ be a field of positive characteristic $p$ dividing $|G|$. Suppose that $\varphi : G \to \text{GL}_n(F)$ is an irreducible representation, show that $\varphi$ is the trivial representation, hence $\chi_\varphi = n1_F$. In particular, if $V$ and $V'$ are finitely generated $F[G]$-modules (hence finite dimensional vector spaces over $F$), then $\chi_V = \chi_{V'}$ if $\dim V \equiv \dim V' \mod p$.

110. Orthogonality Relations

Throughout this section, we shall let $\overline{F}$ denote an algebraic closure of the field $F$.

Definition 110.1. A group $G$ is called a torsion group if every element of $G$ has finite order. (G may be infinite.)

Proposition 110.2. Let $G$ be a torsion group, $V$ an $n$-dimension vector space over $F$, $V$ an $F[G]$-module. Let $x$ be an element of $G$ and $N$ the order of the cyclic subgroup $\langle x \rangle$ in $G$. Then $\chi_V(x)$ is a sum of $n$ terms in which each term is an $N$th root of unity in an algebraic closure $\overline{F}$ of $F$.

Proof. Let $N = \dim_F V$ and $\varphi : G \to \text{GL}_N(F)$ afford $V$. Set $\alpha = \varphi(x)$ and $n = |\langle x \rangle|$. Taking a Jordan canonical form of $\alpha$ in $\text{GL}_N(\overline{F})$ we see that trace $\alpha$ is a sum of $N$ elements each of which is an eigenvalue of $\alpha$ (counted with multiplicity). So it suffices to show that every eigenvalue of $\alpha$ is an $n$th root of unity in $\overline{F}$. If $v \in V$ is a nonzero eigenvector of $\alpha$ with eigenvalue $\varepsilon$, then $x^n = e_G$ implies that $\alpha^n = 1$, so $v = \alpha^n v = \varepsilon^n v$. Consequently, $1 = \varepsilon^n$ as needed.

Recall that the integral closure of $Z$ in $K$, $K/\mathbb{Q}$ a field extension, is denoted by $\mathbb{Z}_K$. Then the proposition implies:

Corollary 110.3. Let $G$ be a torsion group and $n$ a positive integer satisfying $x^n = e$ for every element $x$ in $G$. Suppose that $F$ is an algebraically closed field of characteristic zero and $\omega$ a primitive $n$th root of unity in $F$. If $\varphi : G \to \text{GL}_N(F)$ is a representation, then $\chi_\varphi : G \to \mathbb{Z}_{\mathbb{Q}(\omega)}$.

Definition 110.4. Let $\chi : G \to F$ be a character. The kernel of $\chi$ is defined to be the set

$$\ker \chi := \chi^{-1}(\chi(1)) = \{x \in G \mid \chi(x) = \chi(1)\}.$$ 

Proposition 110.5. Let $G$ be a finite group, $F$ a field of characteristic zero, and $\varphi : G \to \text{GL}_n(F)$ a representation. Then $\ker \chi_\varphi = \ker \varphi$. In particular, $\ker \chi_\varphi$ is a normal subgroup of $G$. 


PROOF. Certainly, \( \ker \varphi \subset \ker \chi \varphi \), so we need only show the reverse inclusion. So suppose that \( \chi \varphi(x) = \chi \varphi(1) \). By the proposition, we have

\[
n = \chi \varphi(1) = \varepsilon_1 + \cdots + \varepsilon_n
\]

with each \( \varepsilon \) an \( |G| \)th root of unity over \( F \) hence over \( \mathbb{Q} \) (and viewed in \( \mathbb{C} \)). So \( n \leq |\varepsilon_1| + \cdots + |\varepsilon_n| \leq n \). Check if \( |\varepsilon_1| + \cdots + |\varepsilon_n| = |\varepsilon_1| + \cdots + |\varepsilon_n| \), then \( \varepsilon_i = \varepsilon_j \) for all \( i, j \). It follows that \( n = \varepsilon_1 + \cdots + \varepsilon_n = n \varepsilon_1 \), hence \( \varepsilon_i = \varepsilon_1 \) for all \( i \). Since the \( \varepsilon_i \) are the eigenvalues of \( \varphi(x) \) in an algebraic closure of \( F \) by the previous proof, we must have the characteristic polynomial of \( \varphi(x) \) is \( f_{\varphi(x)} = (t - 1)^n \). Since \( \varphi(y) \) is a root of \( t|G| - 1 \) for all \( y \in G \), i.e., \( \varphi(y)|G| - 1 = 0 \), the matrix \( \varphi(x) \) is a root of the gcd of \( (t - 1)^n \) and of \( t|G| - 1 \). The latter polynomial cannot have multiple roots, as we are working in characteristic zero, so this gcd is \( t - 1 \). It follows that \( \varphi(x) = I \) and \( x \) lies in \( \ker \varphi \). \( \square 

**Theorem 110.6.** Let \( G \) be a finite group and \( \varphi : G \to \text{GL}_n(F) \) and \( \varphi' : G \to \text{GL}_m(F) \) two irreducible representations. Let

\[
a_{ij}, a'_{kl} : G \to F \text{ satisfy } \varphi(x) = (a_{ij}(x)), \varphi'(x) = (a'_{kl}(x))
\]

for all \( x \) in \( G \), i.e., \( a_{ij}, a'_{kl} \) are the coordinate functions of \( \varphi, \varphi' \), respectively. Then:

1. If \( \varphi \) and \( \varphi' \) are not equivalent, then

\[
\sum_G a_{ij}(x)a'_{kl}(x^{-1}) = 0 \text{ for all } i, j, k, l.
\]

2. Suppose that \( \text{char } F \) is zero or does not divide the order of \( G \) and \( F|G| \) is \( F \)-split. Then \( n \) is nonzero in \( F \) and

\[
\sum_G a_{ij}(x)a'_{kl}(x^{-1}) = \delta ij \delta jk \frac{|G|}{n}.
\]

**Proof.** Let \( V \) with (ordered) basis \( B = \{v_1, \ldots, v_n\} \) be afforded by \( \varphi \) and \( V' \) with (ordered) basis \( B' = \{v'_1, \ldots, v'_m\} \) be afforded by \( \varphi' \). Check if \( f : V \to V' \) is an \( F \)-linear transformation, then \( Tf : V \to V' \) given by

\[
Tf = \sum_G \varphi'(x^{-1}) \circ f \circ \varphi(x),
\]

i.e.,

\[
Tf(\varphi(\sigma)v) = Tf(\sigma(v)) = \sigma Tf(v) = \varphi'(\sigma)Tf(v)
\]

for all \( v \) in \( V \) and for all \( \sigma \) in \( G \), is an \( F|G| \)-homomorphism. (Cf. \( \text{Tr}_G \).) For fixed \( i \) and \( l \), define an \( F \)-linear transformation \( f_o \) by

\[
f_o : V \to V' \text{ given by } \begin{cases} v_i \mapsto v_l \\ v_j \mapsto 0 \text{ if } j \neq i. \end{cases}
\]
(1): By assumption $V \not\cong V'$ as $F[G]$-modules. Since $V$ and $V'$ are irreducible $F[G]$-modules, we must have $f_o = 0$. As $[f_o]_{B,B'} = ((\delta_{ri}\delta_{si})_{rs}) = ((f_o)_{rs})$ and

$$(T f_o)_{kj} = \sum_G \sum_{r,s} (\varphi'(x^{-1}))_{kr}(f_o)_{rs}(\varphi(x))_{sj}$$

(*)

$$= \sum_G \sum_{r,s} a'_{kr}(x^{-1})\delta_{ri}\delta_{si}a_{sj}(x)$$

$$= \sum_G a'_{kl}(x^{-1})a_{ij}(x),$$

we have proven (1) as $T f_o = 0$.

(2): As $V$ is $F[G]$-irreducible and $F$-split, we have $F = \text{End}_{F[G]}(V)$. Let $B' = B$ and $v'_i = v_i$ for all $v_i$, then $T f_o = \lambda(f_o)I$ for some $\lambda(f_o)$ in $F$. Consequently, $\text{trace}(T f_o) = n\lambda(f_o)$. However, we also have

$$\text{trace}(T f_o) = \text{trace}(\sum_G \varphi(x^{-1})[f_o]_{B}\varphi(x))$$

$$= \sum_G \text{trace}(\varphi(x^{-1})[f_o]_{B}\varphi(x))$$

$$= |G| \text{trace}([f_o]_B) = |G| \delta_{il}.$$ 

Hence $n$ is nonzero in $F$ and

$$T f_o = \lambda(f_o)I = \frac{|G|}{n} \delta_{il} I.$$ 

Plugging this into (*) yields

$$\frac{|G|}{n} \delta_{il}\delta_{kj} = (T f_o)_{kj} = \sum_G a_{kl}(x^{-1})a_{ij}(x).$$

□

**Corollary 110.7.** (Frobenius-Schur Theorem) Let $G$ be a finite group and $F$ a field of characteristic zero or with $\text{char}(F) \not| |G|$. Suppose that $F[G]$ is $F$-split and $\varphi^{(k)} : G \to \text{GL}_{n_k}(F)$, $k = 1, \ldots, s$, is a full set of inequivalent irreducible representations of $G$. If $\varphi^{(k)}(g) = (\alpha_{ij}^{(k)}(g))$ for all $g \in G$ and $k = 1, \ldots, s$, then we have:

1. $B := \{a_{ij}^{(k)} \mid i, j = 1, \ldots, n_k, k = 1, \ldots, s\}$ is a linearly independent set.
2. $s = r$, the number of conjugacy classes of $G$.
4. $B$ is a basis for the dual space of $F[G]$, $\text{Hom}_F(F[G], F)$.

**Proof.** Statement (1) follows from Theorem 110.6, (2) and (3) by Summary 107.8, and (4) from (1), (2), and (3). □
Of course, we call two characters $\chi, \chi' : G \to F$ inequivalent if they are not equivalent.

**Theorem 110.8.** (Orthogonal Relations) Let $G$ be a finite group and $\chi, \chi'$ two inequivalent characters.

1. If both $\chi$ and $\chi'$ are irreducible, then
   $$\sum_G \chi(x)\chi'(x^{-1}) = 0.$$ 

2. Suppose that char $F$ is zero or does not divide the order of $G$ and $F[G]$ is $F$-split. If $\chi$ is irreducible of degree $n$, then $\chi(1) = n \neq 0$ in $F$ and
   $$\frac{1}{|G|} \sum_G \chi(x)\chi(x^{-1}) = 1.$$ 

3. If char $F = 0$, then
   a. $\sum_G \chi(x)\chi(x^{-1})$ is an integer.
   b. $|G|$ divides $\sum_G \chi(x)\chi(x^{-1})$ in $\mathbb{Z}$.
   c. If $|G| = \sum_G \chi(x)\chi(x^{-1})$, then $\chi$ is irreducible.

**Proof.** Let $\chi = \chi_\varphi$ and $\chi' = \chi_{\varphi'}$ with $\varphi = (a_{ij})$ and $\varphi' = (a'_{ij})$ the representations giving $\chi$ and $\chi'$ respectively.

1. By the theorem,
   $$\sum_G \chi(x)\chi'(x^{-1}) = \sum_G \left( (\sum_i a_{ii}(x))(\sum_j a_{jj}(x^{-1})) \right) = \sum_{i,j} a_{ii}(x)a'_{jj}(x^{-1}) = 0.$$ 

2. If $n = \chi(1)$, we know $n$ is nonzero in $F$ by the theorem and
   $$\sum_G \chi(x)\chi(x^{-1}) = \sum_G \left( (\sum_i a_{ii}(x))(\sum_j a_{jj}(x^{-1})) \right) = \sum_{i,j} a_{ii}(x)a'_{jj}(x^{-1}) = \sum_{ij} \frac{|G|}{n} = n \frac{|G|}{n} = |G|.$$ 

3. If $\tilde{F}$ is an algebraic closure of $F$, then we can view $\chi$ as a character $\chi : G \to \tilde{F}$. Let $\chi_1, \ldots, \chi_r$ be all the irreducible characters $G \to \tilde{F}$. 

Then \( \chi = \sum m_i \chi_i \) for some non-negative integers \( m_i \). We see that

\[
\sum_G \chi(x) \chi(x^{-1}) = \sum_G \left( \sum_i m_i \chi_i(x) \right) \left( \sum_j m_j \chi_j(x^{-1}) \right)
\]

by (1) and (2). Since \( \sum m_i^2 \) is an integer and characteristic of \( F \) is zero, we have \( \sum_G \chi(x) \chi(x^{-1}) \) dividing \( |G| \).

Finally suppose that \( |G| = \sum_G \chi(x) \chi(x^{-1}) \). Then by (*), we have \( 1 = \sum m_i^2 \), so there exists an \( i \) such that \( m_i = 1 \) and \( m_j = 0 \) for all \( j \neq i \), i.e., \( \chi \) is irreducible as a character when viewed as \( \chi : G \to \tilde{F} \).

**Claim.** \( \chi : G \to F \) is an irreducible:

Let \( V \) be the \( F \)-vector space with basis \( B \) afforded by \( \varphi \) (recall \( \chi = \chi_\varphi \)) and \( W \) a nonzero \( F[G] \)-submodule of \( V \) with \( F \)-basis \( C \). Let \( \tilde{W} \subset \tilde{V} \) be the \( \tilde{F} \)-vector spaces with bases \( \tilde{C} \) and \( \tilde{B} \), respectively. Both are \( \tilde{F}[G] \)-modules with \( \tilde{W} \) a submodule of \( \tilde{V} \). As \( \chi : G \to \tilde{F} \) is afforded by \( \tilde{V} \) and is irreducible, we see that \( \tilde{W} = \tilde{V} \). So

\[
\dim_F W = \dim_{\tilde{F}} \tilde{W} = \dim_{\tilde{F}} \tilde{V} = \dim_F V < \infty.
\]

It follows that \( W = V \) and \( \chi : G \to F \) is irreducible. \( \square \)

**Theorem 110.9.** Let \( G \) be a finite group. Suppose that \( \text{char} F \) is zero or does not divide the order of \( G \) and \( F[G] \) is \( F \)-split. If \( \chi_1, \ldots, \chi_r \) are all the irreducible characters of \( G \) and \( x, y \) lie in \( G \), then

\[
\sum_{i=1}^{r} \chi_i(x) \chi_i(y^{-1}) = \begin{cases} 
0, & \text{if } y \notin C(x), \\
|Z_G(x)|, & \text{if } y \in C(x).
\end{cases}
\]

**Proof.** Let \( x_1, \ldots, x_r' \) be a system of representatives for the conjugacy classes of \( G \). As \( F[G] \) is \( F \)-split, \( r = r' \). Let \( C = (c_{ij}) \) with \( c_{ij} = \chi_i(x_j) \) (independent of the choice of \( x \in C(x_j) \)). This matrix is called the character table of \( G \). Let \( D = (d_{ij}) \) with

\[
d_{ij} = \frac{|C(x_i)|}{|G|} \chi_j(x_i^{-1}) = \frac{1}{|Z_G(x_i)|} \chi_j(x_i^{-1}).
\]
Check that
\[(CD)_{ij} = \sum_k \frac{|C(x_k)|}{|G|} \chi_i(x_k)\chi_j(x_k^{-1}) \]
\[= \frac{1}{|G|} \sum_k \chi_i(x)\chi_j(x^{-1}) \]
\[= \frac{1}{|G|} |G| \delta_{ij} = \delta_{ij}.\]

Thus \(CD = I\), hence \(DC = I\) and
\[(DC)_{ij} = \sum_{k=1}^r \frac{1}{|Z_G(x_i)|} \chi_k(x_i^{-1})\chi_k(x_j).\]

Let \(x \in \mathcal{C}(x_i)\) and \(y \in \mathcal{C}(x_j)\). If \(i \neq j\), we see that \((DC)_{ij} = 0\) leads to
\[\sum_{k=1}^r \chi_k(x)\chi_k(y^{-1}) = 0\]
and if \(i = j\), then \((DC)_{ii} = 1\) leads to
\[\left|Z_G(x_i)\right| = \frac{|G|}{|\mathcal{C}(x)|} \sum_{k=1}^r \chi_k(x)\chi_k(x^{-1}) = \sum_{k=1}^r \chi_k(x)\chi_k(x^{-1}) = \right|Z_G(x_i)| \delta_{ij}.\]

Conclusion 110.10. We have shown if \(\text{char} \ F\) is zero or does not divide the order of \(G\) with \(F[G]\) being \(F\)-split, \(\chi_1, \ldots, \chi_r\) being all the irreducible characters of \(G\), and \(x_1, \ldots, x_r\) being a system of representatives of the conjugacy classes of \(G\), then for all \(i, j = 1, \ldots, r\), and \(x \in G\), we have
\[\sum_{g \in G} \chi_i(x)\chi_j(x^{-1}) = |G| \delta_{ij}\]
\[\sum_{k=1}^r \chi_k(x_i)\chi_k(x_j^{-1}) = \left|Z_G(x_i)\right| \delta_{ij}.\]

Definition 110.11. Let \(F\) be a field and \(G\) a finite group. A function \(f : G \to F\) is called a class function on \(G\) if \(f(g) = f(xgx^{-1})\) for all \(x \in G\). For example, every character of \(G\) is a class function on \(G\). Let \(\text{class}_F(G)\) denote the set of class functions \(G \to F\). This is clearly an \(F\)-vector space. Moreover, \(\text{class}_F(G)\) is an \(r\)-dimensional \(F\)-vector space. Indeed if \(\{C_1, \ldots, C_r\}\) is the set conjugacy classes of \(G\), then the characteristic functions \(f_i \in \text{class}_F(G)\), \(i = 1, \ldots, r\), i.e., \(f_i(x) = \delta_{ij}\) for \(x \in C_j\), form a basis for \(\text{class}_F(G)\).

Corollary 110.12. Let \(G\) be a finite group. Suppose that \(\text{char} \ F\) is zero or does not divide the order of \(G\) and \(F[G]\) is \(F\)-split. Let \(\chi_1, \ldots, \chi_s\) be all the irreducible characters of \(G\). Then \(\{\chi_1, \ldots, \chi_s\}\) is a basis for \(\text{class}_F(G)\).
Proof. The irreducible characters $\chi_1, \ldots, \chi_s$ are linearly independent by the Orthogonality Relations (Theorem 110.8). As $s$ is the number of conjugacy classes by Summary 107.8, the result follows. □

Remark 110.13. Let $G$ be a finite group. We look at complex characters for $G$. Define an inner product on $(\mathbb{C}[G])^\ast = \text{Hom}_\mathbb{C}(\mathbb{C}[G], \mathbb{C})$, the dual space of $\mathbb{C}[G]$, by

$$\langle \chi, \chi' \rangle = \langle \chi, \chi' \rangle_\mathbb{C} : (\mathbb{C}[G])^\ast \times (\mathbb{C}[G])^\ast \to \mathbb{C}$$

given by

$$\langle f, h \rangle = \langle f, h \rangle_\mathbb{C} = \frac{1}{|G|} \sum_G f(x) \overline{h(x)}$$

where $\overline{h(x)}$ is defined by $\overline{h(x)} := \overline{h(x)}$ for $x$ in $\mathbb{C}[G]$. Since $\chi(x)$ is a sum of roots of unity for a character $\chi$ of $G$ and $\chi(x^{-1}) = \overline{\chi(x)}$, we have

$$\langle \chi, \chi' \rangle = \frac{1}{|G|} \sum_G \chi(x) \overline{\chi'(x)} = \frac{1}{|G|} \sum_G \chi(x) \chi'(x^{-1})$$

for characters $\chi, \chi'$ of $G$, where we restrict $\langle \chi, \chi' \rangle$ to the subspace spanned by the characters. We also have

$$\overline{\chi} : G \to \mathbb{C} \text{ given by } \overline{\chi(x)} := \chi(x^{-1})$$

is a character. Moreover, $\chi$ is irreducible if and only if $\overline{\chi}$ is. So complex characters come in two flavors, real, i.e., those have values in $\mathbb{R}$ and complex pairs $\chi, \overline{\chi}$, i.e., those complex characters arising with non-real values.

Let $B = \{\chi_1, \ldots, \chi_r\}$ be the set of all irreducible complex characters of $G$. By Conclusion 110.12 and the Orthogonal Relations (Theorem 110.8), $B$ is an orthonormal basis for the complex inner product space $\mathbb{C}[G]^\ast$ via $\langle \chi, \chi' \rangle$ and if $\chi$ is a complex character of $G$, then $\chi = \sum_{i=1}^r a_i \chi_i$ for some non-negative integers $a_i = \langle \chi, \chi_i \rangle$, $i = 1, \ldots, r$. Therefore, $\langle \chi, \chi \rangle = \sum_{i=1}^r a_i^2 = |G|^2$. It follows by Summary 107.8, if $W$ and $W'$ are finite dimensional $\mathbb{C}[G]$-modules, then

$$\dim_\mathbb{C} \left( \text{Hom}_\mathbb{C}[G](W, W') \right) = \langle \chi_W, \chi_{W'} \rangle_\mathbb{C}.$$

Examples 110.14. 1. Let $G$ be the cyclic group $\langle a \rangle$ of order 3 and $\omega$ a cube root of unity. Then the character table over the complex numbers is:
Let $G$ be the dihedral group $D_3 = \langle a, b \mid a^3 = 1 = b^2, \ bab^{-1} = a^{-1} \rangle$ with $1, a, b$ a system of representatives of the conjugacy classes. Then the character table over the complex numbers is:

<table>
<thead>
<tr>
<th></th>
<th>1</th>
<th>$a$</th>
<th>$a^2$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\chi_1$</td>
<td>1</td>
<td>1</td>
<td>1</td>
</tr>
<tr>
<td>$\chi_2$</td>
<td>1</td>
<td>$\omega$</td>
<td>$\omega^2$</td>
</tr>
<tr>
<td>$\chi_3$</td>
<td>1</td>
<td>$\omega^2$</td>
<td>$\omega$</td>
</tr>
</tbody>
</table>

2. Let $G$ be a finite group. Suppose that char $F$ is zero or does not divide the order of $G$ and $F[G]$ is $F$-split. Let $\mathfrak{A}$ be an irreducible $F[G]$-module and $B_3$ the simple component corresponding to $\mathfrak{A}$. Then show the unit $f = 1_{B_3}$ of $B_3$ satisfies

$$f = \frac{\chi_3(1)}{|G|} \sum_{g \in G} \chi_3(g^{-1}) g.$$

**Exercises 110.15.**

1. (Idempotent Theorem) Let $G$ be a finite group. Suppose that char $F$ is zero or does not divide the order of $G$ and $F[G]$ is $F$-split. Let $\mathfrak{A}$ be an irreducible $F[G]$-module and $B_3$ the simple component corresponding to $\mathfrak{A}$. Then show the unit $f = 1_{B_3}$ of $B_3$ satisfies

(i) $\varphi(g)$ is diagonalizable.

(ii) $\varphi(g)$ is equal to the sum (counted with multiplicity) of the eigenvalues of $\varphi(g)$ all of which are $\chi_V(1)$ $n$th roots of unity.

(iii) $|\chi_V(g)| \leq \chi_V(1)$.

3. Show the complex character $\chi_{C[G]}$ of the regular representation of a finite group $G$ (called the regular character of $G$) satisfies

$$\chi_{C[G]}(x) = \begin{cases} |G| & \text{if } x = 1 \\ 0 & \text{if } x \neq 1. \end{cases}$$
Show if $\chi_V$ is an irreducible complex character of $G$, then $\langle \chi_V, \chi_{C[G]} \rangle = \deg \chi_V$.

4. Let $G$ be a finite group and $N$ a normal subgroup of $G$. Let $\varphi : G \to G/N$ be the canonical epimorphism. Let $\chi$ be a complex character on $G/N$. Show that $\bar{\chi} = \chi \circ \varphi$ is a complex character on $G$ and irreducible if and only if $\chi$ is irreducible. In addition, show that $\bar{\chi}_{C[G/N]} = \sum \chi(1) \chi$ where the sum is taken over all irreducible complex characters on $G$ satisfying $N \subset \ker \chi$ and $\chi_{C[G/N]}$ is the regular character of $G/N$. (Cf. the previous exercise).

5. Let $G$ be the dihedral group $D_4 = \langle a, b \mid a^4 = 1 = b^2, bab^{-1} = a^{-1} \rangle$ with $1, a^2, a, b, ab$ a system of representatives of the conjugacy classes. Show the character table over the complex numbers is

<table>
<thead>
<tr>
<th></th>
<th>1</th>
<th>$a^2$</th>
<th>$a$</th>
<th>$b$</th>
<th>$ab$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\chi_1$</td>
<td>1</td>
<td>1</td>
<td>1</td>
<td>1</td>
<td>1</td>
</tr>
<tr>
<td>$\chi_2$</td>
<td>1</td>
<td>1</td>
<td>1</td>
<td>$-1$</td>
<td>$-1$</td>
</tr>
<tr>
<td>$\chi_3$</td>
<td>1</td>
<td>1</td>
<td>$-1$</td>
<td>1</td>
<td>$-1$</td>
</tr>
<tr>
<td>$\chi_4$</td>
<td>1</td>
<td>1</td>
<td>$-1$</td>
<td>$-1$</td>
<td>1</td>
</tr>
<tr>
<td>$\chi_5$</td>
<td>2</td>
<td>$-2$</td>
<td>0</td>
<td>0</td>
<td>0</td>
</tr>
</tbody>
</table>

6. Fill in the details to Remark 110.13

111. **Burnside’s $p^aq^b$ Theorem**

Throughout this section, we let $\bar{F}$ denote an algebraic closure of the field $F$. Recall that if $K$ is an algebraically closed field of characteristic zero, we shall view an algebraic closure of $\mathbb{Q}$ in $K$ to lie in $\mathbb{C}$ by taking an appropriate isomorphic image.

**Theorem 111.1.** (Arithmetic Lemma) (Frobenius) Let $F$ be a field of characteristic zero with algebraic closure $\bar{F}$ and $G$ a finite group such that $F[G]$ is $F$-split. If $\chi : G \to F$ is an irreducible character and $x$ an element in $G$, then

$$\frac{|C(x)|}{\chi(1)} \chi(x)$$

is an algebraic integer, i.e., lies in $\mathbb{Z}_{\bar{Q}}$, where $\bar{Q}$ is the algebraic closure of $\mathbb{Q}$ in $\bar{F}$.

**Proof.** Let $\chi = \chi_V$ with $V$ the irreducible $F[G]$-module afforded by $\varphi : F[G] \to \text{End}_F(V)$ and $C = C_x = \sum_{C(x)} y$, the class sum of
Let $x$ in $F[G]$. We know that $C$ lies in the center $Z(F[G])$ by Lemma 106.7. By definition, if $z \in Z(F[G])$, then $\varphi(y)\varphi(z) = \varphi(z)\varphi(y)$ for all $y \in F[G]$ means that $\varphi(z)$ lies in $\text{End}_{F[G]}(V)$. As $F[G]$ is $F$-split, $\text{End}_{F[G]}(V) = F1_V$ by Burnside’s Theorem 107.4. Therefore, $\varphi(z)$ lies in $F1_V$ for all $z \in Z(F[G])$. In particular, $\varphi(C) = \lambda 1_V$ for some $\lambda$ in $F$. Let $n = \chi(1) = \dim V$ and $B$ a (ordered) basis for $V$. We know that $\chi(x)$ lies in $Z_{\bar{Q}}$ with $\bar{Q}$ an algebraic closure of $Q$ by Lemma 110.3, so

$$\chi(1)\lambda = n\lambda = \text{trace} \left( [\varphi(C)]_B \right) = \text{trace} \left( [\varphi\left( \sum_{C(x)}(y)) \right)]_B \right) = |C(x)|\chi(x)$$

lies in $Z_{\bar{Q}}$, hence

$$\lambda = \frac{|C(x)|}{\chi(1)}\chi(x) \text{ lies in } \bar{Q}.$$ 

We must show that $\lambda$ lies in $Z_{\bar{Q}}$.

We restrict the map $\varphi$ to a ring homomorphism (hence a $Z$-algebra homomorphism)

$$\varphi = \varphi|_{Z[G]} : Z[G] \rightarrow \text{End}_F(V).$$

Since $C = \sum_{C(x)} 1 \cdot z$ lies in $Z(G)$, it must lie in $Z(Z[G])$. Further restrict $\varphi$ to a ring homomorphism

$$\varphi = \varphi|_{Z(Z[G])} : Z(Z[G]) \rightarrow \text{End}_F(V).$$

As above, we have $\varphi\left( Z(Z[G]) \right) \subset \text{End}_{F[G]}(V) = F1_V$, i.e., this restriction means $\varphi : Z(Z[G]) \rightarrow F1_V = \text{End}_{F[G]}(V)$ and $C$ lies in $Z(Z[G])$. Since $\bar{Q}$ is a finite group, $Z[G]$ is a finitely generated abelian group, hence so is the subring $\varphi\left( Z(Z[G]) \right)$ – using either $Z(Z[G])$ is generated by class sums or $Z$ is noetherian. In particular, every element of $\varphi\left( Z(Z[G]) \right)$ is integral over $Z$ by Claim 76.4. So $\varphi(C) = \lambda 1_V$ implies that $\lambda$ lies in $\bar{Q} \cap Z_F \subset \bar{Q} \cap Z_{\bar{F}} = Z_{\bar{Q}}$, where $\bar{F}$ is an algebraic closure of $F$. \hfill \Box

**Theorem 111.2.** Let $F$ be a field of characteristic zero and $G$ a finite group such that $F[G]$ is $F$-split. If $\chi : G \rightarrow F$ is an irreducible character, then $\chi(1) \mid |G|$ in $\mathbb{Z}$.

**Proof.** Let $x_1, \ldots, x_r$ be a system of representatives for the conjugacy classes of $G$ and $\bar{Q}$ an algebraic closure of $Q$. Then

$$\frac{|C(x_i)|}{\chi(1)} \chi(x_i)$$
and \( \chi(x^{-1}) \) lie in \( \mathbb{Z}_{\tilde{Q}} \), hence
\[
\frac{|G|}{\chi(1)} = \sum_{G} \frac{1}{\chi(1)} \chi(x) \chi(x^{-1}) = \sum_{i=1}^{r} \frac{C(x_i)}{\chi(1)} \chi(x) \chi(x^{-1})
\]
lies in \( \mathbb{Z}_{\tilde{Q}} \cap \mathbb{Q} = \mathbb{Z} \). □

Schur proved a stronger result, viz., under the hypothesis of the theorem, \( \chi(1) | [G : Z(G)] \) in \( \mathbb{Z} \). We shall prove this in the Addendum 112.

**Corollary 111.3.** Let \( F \) be a field of characteristic zero and \( G \) a finite group such that \( F[G] \) is \( F \)-split, say \( F[G] \cong \times_{i=1}^{r} M_{n_i}(F) \). Then

1. \( r \) is the number of conjugacy classes of \( G \).
2. \( |G| = \sum_{i=1}^{r} n_i^2 \).
3. There exists an \( i, 1 \leq i \leq r \) satisfying \( n_i = 1 \).
4. \( n_i | |G| \) for \( i = 1, \ldots, r \).

**Proof.** We have previously shown (1) — (3) and (4) follows from the theorem. □

**Lemma 111.4.** (Burnside’s Lemma) Let \( F \) be an algebraically closed field of characteristic zero, \( G \) a finite group, and \( \varphi : G \to \text{GL}_n(F) \) an irreducible representation. If \( x \) in \( G \) satisfies \(|C(x)| \) is relatively prime to \( n = \chi_{\varphi}(1) \), then either \( \chi_{\varphi}(x) = 0 \) or \( \varphi(x) = \lambda I \) for some \( \lambda \) in \( \mathbb{Z}_{\tilde{Q}} \), with \( \tilde{Q} \) the algebraic closure of \( \mathbb{Q} \).

**Proof.** Let \( \chi = \chi_{\varphi} \) and \( m = |C(x)| \), so \( m \) and \( n \) are relatively prime. We can write \( 1 = am + bn \) for some integers \( a, b \); hence
\[
\lambda := \frac{\chi(x)}{n} = a \frac{m}{n} \chi(b) + b \chi(x)
\]
lies in \( \mathbb{Z}_{\tilde{Q}} \), with \( \tilde{Q} \) an algebraic closure of \( \mathbb{Q} \) by the Arithmetic Lemma. By Proposition 110.2 and the proof of Proposition 110.5, we know that \( \chi(x) = \varepsilon_1 + \cdots + \varepsilon_n \) for some \( |G| \) th roots of unity \( \varepsilon_i \) in \( \tilde{Q} \) that are also the eigenvalues of \( \varphi(x) \) and satisfy
\[
|\lambda| = \frac{1}{n} |\chi(x)| = \frac{1}{n} |\varepsilon_1 + \cdots + \varepsilon_n| \leq \frac{1}{n} \sum_{i=1}^{n} |\varepsilon_i|
\]
with equality if and only if \( \varepsilon_i = \varepsilon_j \) for all \( i, j \).

**Case 1.** Not all the \( \varepsilon_i \) are equal:

By \((*)\), we have \( |\lambda| < 1 \). Let \( \omega \) be a primitive \( |G| \) th root of unity over \( \mathbb{Q} \) and \( L = \mathbb{Q}(\omega) \). Then \( \lambda \) lies in \( \mathbb{Z}_{\tilde{Q}} \cap \mathbb{Q} = \mathbb{Z}_L \) and \( \varepsilon_k \) lies in \( \mathbb{Z}_L \) for all \( k \). Suppose that \( \varepsilon_i \neq \varepsilon_j \). Since \( L/\mathbb{Q} \) is a finite Galois
Theorem 111.5. (Burnside) Let $G$ be a finite group and $p$ a prime. If there exists an element $x$ in $G$ satisfying $C(x) = p^r > 1$, then $G$ is not a simple group.

**Proof.** Let $\chi_1, \ldots, \chi_r : G \to \mathbb{C}$ be all the distinct irreducible complex characters. One of these must be the trivial character, say it is $\chi_1$. Let $n_i = \chi_i(1)$, so $n_1 = 1$ (and $\chi_1(x) = 1$). Since $C(x) \neq \{x\}$, we know that $x \neq 1$, hence $1 \notin C(x)$. Therefore, we have

$$0 = \sum_{i=1}^{r} \chi_i(x)\chi_i(1^{-1}) = \sum_{i=1}^{r} \chi_i(1)\chi_i(x) = 1 + \sum_{i=2}^{r} n_i\chi_i(x).$$

**Case 1.** There exists an integer $1 < j \leq r$ such that $p \nmid n_j$ and $\chi_j(x) \neq 0$:

By assumption, $\chi_j(1)$ and $|C(x)|$ are relatively prime. Since $\chi_j(x) \neq 0$, by Burnside’s Lemma, there exists a complex number $\lambda$ satisfying $\varphi_j(x) = \lambda x$, where $\chi_j = \chi_{\varphi_j}$. In particular, $\varphi_j(x)$ lies in $Z(M_{n_j}(\mathbb{C}))$, so $\varphi_j(x)\varphi_j(y) = \varphi_j(y)\varphi_j(x)$ for all $y \in G$. As $|C(x)| > 1$, $x \notin Z(G)$, so there exists a $y \in G$ satisfying $1 \neq xyx^{-1}y^{-1}$ lying in $\ker\varphi_j$. As $j > 1$, $\chi_j$ is not the trivial character, so $1 < \ker\varphi_j = \ker\chi_j < G$ by Proposition 110.5. Consequently, $G$ is not simple.

**Case 2.** For all $1 < j \leq r$, we have $\chi_j(x) = 0$ whenever $p \nmid n_j$:

In this case, we have

\[
(*) \quad 0 = 1 + \sum_{i=2}^{r} n_j\chi_j(x) = 1 + p \sum_{p > 1 \atop p|n_j} \frac{n_j}{p} \chi_j(x).
\]

As $n_j/p$ is an integer when $p|n_j$, we have $1/p$ lies in $\mathbb{Z}_C \cap \mathbb{Q} = \mathbb{Z}$ by $(*)$, a contradiction. Thus Case 2 cannot occur. \qed
Theorem 111.6. (Burnside’s $p^aq^b$-Theorem) Let $G$ be a finite group of order $p^aq^b$ with $p, q$ distinct primes and $a, b$ non-negative integers. Then $G$ is a solvable group.

Proof. By our previous work and induction, it suffices to show that $G$ is not simple if $a$ and $b$ are both positive. Let $Q$ be a Sylow $q$-subgroup of $G$ and $e_G \neq x$ an element in $Z(Q)$. Then $Q \subset Z_G(x)$, hence $|C(x)| = [G : Z_G(x)] = p^n$ for some integer $n \geq 0$. If $n = 0$, then $Z_G(x) = G$, so $1 < Z(G) \triangleleft G$ as $x \neq e_G$, and the result follows. If $n > 0$, then $G$ is not simple by Burnside’s Theorem. □

When Burnside wrote the first edition of his historic book on finite group theory, he decided not discuss representation theory as it was not intrinsic to the theory of groups. After he proved the $p^aq^b$-Theorem in 1904, he realized that representation theory was now an essential tool in studying group theory, so he included it in a revised edition of his book. A proof of his theorem was proven avoiding representation in the 1970’s. It is much more difficult than the one using representation theory.

112. Addendum: Schur’s Theorem

In this section, we establish the improvement to Theorem 111.2 mentioned in §111. To do so, we assume the reader is conversant with tensor products. (Cf. §116.)

Definition 112.1. Let $F$ be a field and $G$ a group. Suppose that $V$ is an irreducible $F[G]$-module that is finite dimensional as an $F$-vector space. We say that $V$ is absolutely irreducible if $\tilde{F} \otimes_F V$ is an irreducible $\tilde{F}[G]$-module with $\tilde{F}$ an algebraic closure of $F$.

Check 112.2. An $F[G]$-module $V$, finite dimensional as an $F$-vector space, is absolutely irreducible if and only if $\text{End}_{F[G]}(V) = F$.

Lemma 112.3. Let $G, G'$ be (arbitrary) groups. Suppose that $V$ is an absolutely irreducible $F[G]$-module and $V'$ is an absolutely irreducible $F[G']$-module. Let $V \otimes_F V'$ be the $F[G \times G']$-module induced by the $G$-action $(g, g')(v \otimes v') = gv \otimes g'v'$ for all $g \in G, g' \in G', v \in V,$ and $v' \in V'$. Then $V \otimes_F V'$ is an absolutely irreducible $F[G \times G']$-module.

Proof. Let $\varphi : F[G] \to \text{End}_F(V)$ and $\varphi' : G' \to \text{End}_F(V')$ afford $V$ and $V'$, respectively. By (the modified form Proposition 107.6 of) Burnside’s Theorem 107.4, we know both $\varphi$ and $\varphi'$ are surjective.
Consider the following commutative diagram:

\[
\begin{array}{ccc}
F[G \times G'] & \xrightarrow{\rho} & \text{End}_F(V \otimes_F V') \\
\downarrow f & & \uparrow g \\
F[G] \otimes_F F[G'] & \xrightarrow{h} & \text{End}_F(V \otimes_F V'),
\end{array}
\]

where

\[
\rho: (g, g') \mapsto (v \otimes v' \mapsto gv \otimes g'v')
\]

\[
f: (g, g') \mapsto g \otimes g'
\]

\[
g: h \otimes h' \mapsto (v \otimes v' \mapsto h(v) \otimes h'(v')),
\]

for all \(v \in V, v' \in V', g \in G, g' \in G', h \in \text{End}_F(V), \) and \(h' \in \text{End}_F(V')\), induce the \(F\)-algebra homomorphisms in the diagram. By dimension count, the \(F\)-algebra homomorphism \(g\) is an isomorphism (as it is \(F\)-linear). The map \(f\) is clearly a surjection. As the map \(F[G] \otimes_F F[G'] \rightarrow F[G \times G']\) induced by \(g \otimes g' \mapsto (g, g')\) determines the inverse to \(f\), we conclude that \(f\) is also an \(F\)-algebra isomorphism.

Since \(\varphi \otimes \varphi'\) is surjective, it follows that \(\rho\) is also surjective. Taking the fixed points of the action of \(G \times G'\) on \(V \otimes_F V'\), we see that

\[
\text{End}_{F[G \times G']}(V \otimes_F V') = \left(\text{End}_F(V \otimes V')\right)^{G \times G'} = F.
\]

Therefore, by Check 112.2, \(V \otimes_F V'\) is an absolutely irreducible \(F[G \times G']\)-module. \(\Box\)

**Theorem 112.4.** (Schur) Let \(F\) be a field of characteristic zero and \(G\) a finite group. Suppose that \(F[G]\) is \(F\)-split and \(\chi : G \rightarrow F\) is an irreducible character. Then \(\chi(1) \mid [G : Z(G)]\).

**Proof.** (Tate). Let \(\chi = \chi_V\). By Check 112.2, \(V\) is absolutely irreducible as \(\text{End}_F(V) = F\). Let \(n \in \mathbb{Z}^+\). Then \(V^\otimes_n := V \otimes_F \cdots \otimes_F V\) is an absolutely irreducible \(F[G \times \cdots \times V]\)-module by Lemma 112.4. Let the map \(\varphi : G \rightarrow \text{Aut}_F(V)\) afford \(V\). Since \(\text{End}_{F[G]}(V) = F\) and \(\varphi(Z(G)) \subset \text{End}_{F[G]}(V)\), as \(\varphi(Z(G))\) commutes with \(\varphi(G)\), we conclude that \(\varphi|_{Z(G)} : Z(G) \rightarrow F^\times\). Set

\[H = \{(g_1, \ldots, g_n) \mid g_i \in Z(G), g_1 \cdots g_n = 1\} \subset G \times \cdots \times G,\]

a subgroup. For all \((g_1, \ldots, g_n) \in H\) and \(v_1, \ldots, v_n \in V\), we have

\[
(g_1, \ldots, g_n)(v_1 \otimes \cdots \otimes v_n) = g_1v_1 \otimes \cdots \otimes g_nv_n
\]

\[
= \varphi(g_1)v_1 \otimes \cdots \otimes \varphi(g_n)v_n = \varphi(g_1) \cdots \varphi(g_n)v_1 \otimes \cdots \otimes v_n
\]

\[
= \varphi(g_1 \cdots g_n)(v_1 \otimes \cdots \otimes v_n) = v_1 \otimes \cdots \otimes v_n,
\]
as \( \varphi(g_i) \) lies in \( \text{End}_F(V) = F \) for all \( g_i \in Z(G) \), \( i = 1, \ldots, n \). Since the \( v_1 \otimes \cdots \otimes v_n \) generate \( V^{\otimes n} \), it follows that \( H \) acts trivially on \( V^{\otimes n} \).

Since \( H \subset Z[G \times \cdots \times G] \), we have \( H \triangleleft G \times \cdots \times G \). By the above, the \( H \)-action on \( V^{\otimes n} \) is trivial. Consequently, we may view \( V^{\otimes n} \) as an irreducible \( F[G \times \cdots \times G]/H \)-module. By Theorem 111.2, we see that

\[
\chi(1)^n \bigg| \frac{|G|^n}{|Z(G)|^{n-1}}.
\]

Let \( g_1, \ldots, g_{n-1} \in Z(G) \). If \( g_n \in G \), then \( g_1 \cdots g_{n-1} g_n \in H \) if and only if \( g_n = (g_1 \cdots g_{n-1})^{-1} \). Therefore,

\[
\chi(1)^n \bigg| \frac{|G|^n}{|Z(G)|^{n-1}}.
\]

Hence, there exists an element \( e = e(m) \in Z \) satisfying \( |G|^n/|Z(G)|^{n-1} = e\chi(1)^n \). It follows that

\[
\left( \frac{|G|}{|Z(G)|} \right)^n \bigg| \frac{1}{|Z(G)|} \text{ for all } n \in \mathbb{Z}^+. \text{ Since } [G : Z(G)]/\chi(1) \text{ lies in } \mathbb{Q}, \text{ it follows that we have } [G : Z(G)]/\chi(1) \text{ lies in } \mathbb{Z}, \text{ i.e., } \chi(1) \bigg| [G : Z(G)].
\]

\[ \square \]

**Exercises 112.6.**

1. Prove Check 112.2.

2.

3. Show an \( F[G] \)-module \( V \), finite dimensional as an \( F \)-vector space, is absolutely irreducible if and only if \( K \otimes_F V \) is irreducible for all field extensions \( K/F \).

4. Verify Equation 112.5

### 113. Induced Representations

In this section, we shall study how subgroups of a finite group induce representations of the full group as well as the induced character theory.

In this section we shall assume familiarity of tensor products of modules over an arbitrary ring. (Cf. Section 116.)

**Definition 113.1.** Let \( F \) be a field and \( G \) a finite group. If \( V \) is a finite dimensional \( F[G] \)-module and \( H \subset G \) a subgroup, then restriction of scalars to \( F[H] \) induces an \( F[H] \)-module structure of \( V \) called the **restriction** of \( V \) to \( H \) and denoted by \( \text{res}_H^G(V) \). If \( V \) is a finite dimensional \( F \)-vector space and affords the representation \( \sigma_V : G \to \text{GL}(V) \), then we write \( \text{res}_H^G(\sigma_V) \) for the representation afforded by \( \text{res}_H^G(V) \) and \( \text{res}_H^G \chi \) for the character of this representation. If \( W \) is an \( F[H] \)-module,
then define the induced $F[G]$-module by $\text{ind}^G_H(W) := F[G] \otimes_{F[H]} W$. For example by properties of $\otimes$, we have $\text{ind}^G_H(F[H]) = F[G] \otimes_{F[H]} F[H]$ which is identified with $F[G]$. Since $F[G]$ is an $(F[G], F[H])$-bimodule, we have $\text{ind}^G_H(W) = \{ \sum_{g_i \in G} g_i \otimes w_i \mid g_i \in G, \, w \in W \}$ becomes an $F[G]$-module via the $G$-action $g(\sum_{g_i \in G} g_i \otimes w) = \sum_{g_i \in G} g g_i \otimes w$. If $W$ is a finite dimensional $F$-vector space and affords $\sigma_W : H \to \text{GL}(W)$, let $\text{ind}^G_H(\chi_W)$ denote the character induced by $\text{ind}^G_H(W)$.

**Construction 113.2.** Let $F$ be a field, $G$ a finite group and $H \subset G$ is a subgroup. Suppose that $g_1, \ldots, g_n$ is a left transversal of $H$ in $G$ with $g_1 = e_G = 1_{F[G]} = 1$. So we have $F[G] = \bigoplus_{i=1}^m g_i F[H]$. Therefore, if $W$ is an $F[H]$-module,

$$\text{ind}^G_H(W) = F[G] \otimes_{F[H]} W = \bigoplus_{i=1}^n (g_i F[H] \otimes_{F[H]} W) = \bigoplus_{i=1}^n g_i \otimes_{F[H]} W$$

as $\bigoplus$ and $\otimes_{F[H]}$ commute and $\otimes_{F[H]}$ is $F[H]$-balanced. We have $g_i \otimes_{F[H]} W = 1 \otimes_{F[H]} W$ an $F[H]$-submodule of $\text{res}^G_H(\text{ind}^G_H(W))$. So $w \mapsto 1 \otimes w$ for $w \in W$ induces an $F[H]$-monomorphism $W \to \text{ind}^G_H(W)$. We also have that the direct summand $g_i(1 \otimes_{F[H]} W)$ of $\text{ind}^G_H(W)$ can be expressed as $g_i(1 \otimes_{F[H]} W)$. Therefore, we have $g_i(1 \otimes_{F[H]} W) \cong g_i \otimes_{F[H]} W$ as vector spaces over $F$. For each $i$, $i = 1, \ldots, n$, and $g \in G$, there exists a unique $j$, $1 \leq j \leq n$, satisfying $g g_i \in g_j H$. So for each $i$, there also exists a unique $h_i$ in $H$ satisfying $g g_i = g_j h_i$. It follows that $g$ permutes the $\{g_i \otimes_{F[H]} W \}$ with action $g(g_i \otimes w) = g_j h_i \otimes w = g_j \otimes h_i w$ for all $w \in W$.

Now suppose that $W$ is a finite dimensional $F$-vector space. Let $\sigma_W : H \to \text{GL}(W)$ afford $W$ and $\{w_1, \ldots, w_m \}$ be an $F$-basis for $W$. Then for all $h \in H$, we have $hw_j = \sum_{i=1}^m \alpha_{ij}(h) w_i$ with all $\alpha_{ij}(h) \in F$. It follows that $B := \{g_i \otimes w_j \mid 1 \leq i \leq n, \, 1 \leq j \leq m \}$ is an $F$-basis for $\text{ind}^G_H(W)$ as $\bigoplus$ and $\otimes$ commute. With this notation, we have

$$g(g_i \otimes w_s) = g_j h_i \otimes w_s = g_j \otimes h_i w_s = \sum_{r=1}^m \alpha_{rs}(h_i)(g_j \otimes w_r) \text{ for all } 1 \leq s \leq m.$$

Since $g g_i = g_j h_i$ above means that $h_i = g_j^{-1} g g_i$, we can extend the maps $\alpha_{ij} : H \to F$ to maps $\hat{\alpha}_{ij} : G \to F$ defined by

$$\hat{\alpha}_{ij}(g) = \begin{cases} \alpha_{ij}(g) & \text{if } g \in H \\ 0 & \text{if } g \notin H. \end{cases}$$
Then the action of $g$ in $G$ on the $F$-basis $\mathcal{B}$ for $\text{ind}_H^G(W)$ is given by
\[
g(g_i \otimes w_s) = \sum_{r=1}^{m} \alpha_{rs}(h_i)(g_j \otimes w_r) = \sum_{j=1}^{n} \sum_{r=1}^{m} \hat{\alpha}_{rs}(g_j^{-1}gg_i)(g_j \otimes w_r).
\]

Now extend $\sigma_W$ to an $F[G]$-module by $\hat{\sigma}_W(g) = (\hat{\alpha}_{ij}(g))$. Relative to the $F$-basis $\mathcal{C} = \{g_1 \otimes w_1, \ldots, g_1 \otimes w_m, g_2 \otimes w_1, \ldots, g_n \otimes w_m\}$, the matrix representation of $\text{ind}_H^G(W)$ is
\[
g \mapsto \sigma_{\text{ind}_H^G(W)}(g) = \begin{pmatrix} \hat{\sigma}(g_1^{-1}gg_1) & \cdots & \hat{\sigma}(g_1^{-1}gg_m) \\ \vdots & & \vdots \\ \hat{\sigma}(g_n^{-1}gg_1) & \cdots & \hat{\sigma}(g_n^{-1}gg_m) \end{pmatrix}.
\]

It follows that $\text{ind}_H^G(\chi_W)$ is defined by
\[
\text{ind}_H^G(\chi_W)(g) = \sum_{i=1}^{n} \hat{\chi}_W(g_i^{-1}gg_i)
\]
where
\[
\hat{\chi}_W(g) = \begin{cases} \chi(g) & \text{if } g \in H \\ 0 & \text{if } g \notin H. \end{cases}
\]

Now suppose that $|H| \in F^\times$, e.g., $\text{char}(F) = 0$. Then Equation 113.3 becomes
\[
\text{ind}_H^G(\chi_W)(g) = \frac{1}{|H|} \sum_{x \in G} \hat{\chi}_W(x^{-1}gx),
\]
independent of the transversal of $H$ in $G$. Moreover, in this case, we have
\[
\deg \left( \text{ind}_H^G(\chi_W) \right) = \deg \left( \text{ind}_H^G(\chi_W)(1) \right) = [G:H]\chi_W(1).
\]

**Example 113.6.** Let $F$ be a field, $G$ a finite group and $H$ a subgroup of $G$ of index $n$. Let $1_H$ be the trivial representation of $H$ and $W = Fw$ afford $1_W$. So $hw = w$ for all $h \in H$. In the notation of the construction above, the $G$-action of $g \in G$ on the summand $F(g_i \otimes w)$ that affords $\text{ind}_H^G(1_W)$ is $g(g_i \otimes w) = g_i \otimes w$ with $gg_i \in g_iH$. It follows that the $G$-set of $F$-subspaces $\{F(g_i \otimes w)\}$ of $\text{ind}_H^G(\chi_W)$ with this $G$-action is isomorphic to the $G$-set of left cosets $G/H$ with action left translation on the basis $\mathcal{D} = \{w_i \mid i = 1, \ldots, n\}$ with the $w_i$ corresponding to $g_i \otimes w$. In particular, $1_H^G$ is equivalent to the permutation representation of $G$ afforded by the $G$-set $G/H$ by choosing an $F$-basis $\mathcal{D}$ defined by $gw_i = w_j$ if $gg_iH = g_jH$ for $i = 1, \ldots, n$. In particular, the left regular representation $F[G]$ is isomorphic to $\text{ind}_H^G(1_H)$ when $H = 1$. 
Proposition 113.7. Let $F$ be a field and $G$ a finite group, $H \subset G$ a subgroup of index $n$. Let $V$ be an $F[G]$-module so that $\text{res}_G^H(V)$ contains an $F[H]$-submodule $W$ that satisfies $V = \bigoplus_{i=1}^n g_i W$, where $\{g_1, \ldots, g_n\}$ is an transversal of $H$ in $G$. Then $V = \text{ind}_H^G(\chi_W)$ as an $F[H]$-module.

Proof. The map $F[G] \times W \to V$ is checked to be balanced, so induces a surjective $F$-linear transformation $\varphi : F[G] \otimes_{F[H]} W \to V$ satisfying $\varphi(g_i \otimes F[H] W) = g_i W$, $i = 1, \ldots, n$. As $V = \bigoplus_{i=1}^n g_i W$, we have a linear transformation $V \to \text{ind}_H^G(\chi_W)$ inverse to $\varphi$, so $\varphi$ is an $F$-isomorphism. This map is, in fact, an $F[G]$-homomorphism, so an $F[G]$-isomorphism.

We leave the following as an exercise.

Proposition 113.8. Let $F$ be a field, $G$ a finite group and $H_0 \subset H$ subgroups of $G$.

1. If $W_1$ and $W_2$ are $F[H]$-modules, then
   \[ \text{ind}_G^H(W_1 \oplus W_2) \cong \text{ind}_G^H(W_1) \oplus \text{ind}_G^H(W_2). \]

2. If $W_0$ is an $F[H_0]$-module, then
   \[ \text{ind}_G^H(\text{ind}_{H_0}^H(W_0)) \cong \text{ind}_{H_0}^G(W_0). \]

3. If $W_0$ is a finite dimensional $F[H_0]$-module, then
   \[ \text{ind}_G^H(\text{ind}_{H_0}^H(\chi_{W_0})) \cong \text{ind}_{H_0}^G(\chi_{W_0}). \]

Theorem 113.9. (Frobenius Reciprocity – Module Form) Let $F$ be a field and $G$ a finite group, $H \subset G$ a subgroup. If $V$ is an $F[G]$-module and $W$ is an $F[H]$-module, then

\[ \text{Hom}_{F[H]}(W, \text{res}_G^H(V)) \cong \text{Hom}_{F[G]}(\text{ind}_G^H(W), V). \]

Proof. By the Adjoint Associativity Theorem 116.16, we have a group isomorphism

\[ \text{Hom}_{F[H]}(W, \text{Hom}_{F[G]}(F[G], V)) \cong \text{Hom}_{F[G]}(\text{ind}_G^H(W), V) \]

that is easy to see is an isomorphism of $F$-vector spaces. Therefore, it suffices to show that $\text{Hom}_{F[G]}(F[G], V) \cong \text{res}_G^H(V)$ as $F[H]$-modules where $\text{Hom}_{F[G]}(F[G], V)$ is an $F[H]$-module by $(af)(x) := f(ax)$ for all $f \in \text{Hom}_{F[G]}(F[G], V)$, $x \in F[G]$, and $a \in F[H]$. Let $\theta : \text{Hom}_{F[G]}(F[G], V) \to V$ be the linear transformation defined by $f \mapsto f(1)$. It is easy to see that this is an $F$-vector space isomorphism. Then for all $x \in F[G], a \in F[H], f \in \text{Hom}_{F[G]}(F[G], V)$, we have

\[ \theta(af) = (af)(1) = f(a) = a(f(1)) = a\theta(f). \]

It follows that $\theta$ is an $F[H]$-isomorphism.
Remark 113.10. Let $G$ be a finite group and $F$ a subfield of $\mathbb{C}$. Suppose that $F[G]$ is $F$-split. Of course, $G$ is $\mathbb{C}$-split also. Let $K$ be the field composite of $F$ and $\mathbb{C}$ in some larger field (with fixed embeddings). Then $G$ is also $K$-split by extension of scalars. The basic sets of irreducible $F[G]$-modules and basic sets of irreducible $\mathbb{C}[G]$-modules yield isomorphic basic sets of irreducible $K[G]$-modules. Moreover the characters viewed as class functions on each are then all the same. Therefore, we can identify the irreducible characters on $F$ and $\mathbb{C}$. In particular, Remark 110.13 holds for any such $F$ in $\mathbb{C}$.

Therefore, by Remarks 110.13 and 113.10 we have:

Corollary 113.11. Let $F$ be a subfield of the complex numbers and $G$ a finite group, $H \subset G$ a subgroup. Suppose that both $F[G]$ and $F[H]$ are $F$-split. If $V$ is a finite dimensional $F[G]$-module and $W$ is a finite dimensional $F[H]$-module, then

1. $\dim_F \big( \text{Hom}_{F[H]}(W, \text{res}^G_H(V)) \big) = \langle \chi_W, \text{res}^G_H(\chi_V) \rangle_H$.
2. $\dim_F \big( \text{Hom}_{F[G]} \big( \text{ind}^G_H(W), V \big) \big) = \langle \text{ind}^G_H(\chi_W), \chi_V \rangle_G$.

In particular,

$$\langle \chi_W, \text{res}^G_H(\chi_V) \rangle_H = \langle \text{ind}^G_H(\chi_W), \chi_V \rangle_G.$$

Recall the $F$-vector space class $F(G)$ of a finite group $G$ has as an $F$-basis the irreducible characters under the conditions that char $F = 0$ or char $G$ and $F[G]$ is $F$-split by Corollary 110.12. We generalize the character form of Frobenius Reciprocity without using the previous module form.

Theorem 113.12. (Frobenius Reciprocity – Class Function Case) Let $F$ be a subfield of the complex numbers, $G$ a finite group, and $H$ a subgroup of $G$. Let $\lambda \in \text{class}_F(H)$ and $\mu \in \text{class}_F(G)$. Define $\text{ind}^G_H(\lambda) : G \to H$ by

$$\text{ind}^G_H(\lambda) := \frac{1}{|H|} \sum_{x \in G} \hat{\lambda}(x^{-1}gx),$$

where $\hat{\lambda}$ is the function $\lambda$ extended to $G$ (analogous to Equation 113.4). Then $\text{ind}^G_H(\lambda) \in \text{class}_F(G)$ and

$$\langle \text{ind}^G_H(\lambda), \mu \rangle_G = \langle \lambda, \text{res}^G_H(\mu) \rangle_H.$$

Proof. The first statement is clear. As for the equation, since $\mu$ is a class function and for each $y \in G$, there exist $|G|$ $x$’s satisfying $y = x^{-1}gx$. Therefore, we have, using Remark 110.13 and Equation
Examples 113.13. 1. Let $G$ be a finite group and $H \subset G$ a finite subgroup. Suppose that $\lambda$ is a character of $H$ over $\mathbb{C}$. Then the class function $\text{ind}^G_H(\lambda)$ is character of $G$ over $\mathbb{C}$ if and only if $\langle \text{ind}^G_H(\lambda), \chi \rangle_G = \langle \lambda, \text{res}^G_H(\chi) \rangle_H$ is a non-negative integer for all irreducible characters $\chi$ of $G$ over $\mathbb{C}$ as $\text{res}^G_H(\chi)$ is a character of $H$ over $\mathbb{C}$.

2. Let $D_m$ be the dihedral group of order $2m$ with the presentation $\langle a, b \mid a^m = b^2 = (ab)^2 = 1 \rangle$. Let $A = \langle a \rangle$ and $\lambda : A \to \mathbb{C}$ be a linear character (hence irreducible). Let $\zeta$ be a primitive $m$th root of unity. So $\lambda(a) = \zeta^i$ for some $i$. Then $\text{ind}^D_A(m)(x) = \hat{\lambda}(x) + \hat{\lambda}(bxb^{-1})$ for $x \in A$ satisfying $xbx^{-1} = x^{-1}$ by Equation 113.4. Therefore, $\text{res}^D_A(m)\left(\text{ind}^D_A(m)(\lambda)\right) = \lambda + \bar{\lambda}$. By Frobenius Reciprocity, we have

$$\langle \text{ind}^D_A(m)(\lambda), \text{ind}^D_A(m)(\lambda) \rangle_{D_m} = \langle \lambda, \text{res}^D_A(m)\left(\text{ind}^D_A(m)(\lambda)\right) \rangle_A = \langle \lambda, \lambda + \bar{\lambda} \rangle.$$ 

By Orthogonal Relations (Theorem 110.8), $\text{ind}^D_A(m)(\lambda)$ is irreducible if and only if $\langle \text{ind}^D_A(m)(\lambda), \text{ind}^D_A(m)(\lambda) \rangle = 1$. It follows that $\text{ind}^D_A(m)(\lambda)$ is irreducible if and only if $\lambda \neq \bar{\lambda}$. If $\mu : A \to \mathbb{C}$ is a character with $\text{ind}^D_A(m)(\lambda)$ and $\text{ind}^D_A(m)(\mu)$ both irreducible, then

$$\langle \text{ind}^D_A(m)(\lambda), \text{ind}^D_A(m)(\mu) \rangle_{D_m} = \langle \lambda, \mu + \bar{\mu} \rangle_A = \begin{cases} 0 & \text{if } \lambda \neq \mu \text{ and } \lambda \neq \bar{\mu} \\ 1 & \text{if } \lambda = \mu \text{ or } \lambda = \bar{\mu}. \end{cases}$$

There are two possible cases.

**Case 1.** $m$ is odd:

In this case, there are two linear characters of $D_m$ over $\mathbb{C}$ as the commutator of $D_m$ has index two and $\frac{1}{2}(m-1)$ distinct characters of the form $\text{ind}^D_A(m)(\lambda)$ by the argument above. Therefore, these irreducible characters constitute $1^2 + 1^2 + 2^2 \left(\frac{1}{2}(m-1)\right) = 2 + 2(m-1) = 2m$ in the sum of squares giving $|D_m|$. It follows that these must be all the irreducible characters.

**Case 2.** $m$ is even:
In this case, there are four linear characters of $D_m$ over $\mathbb{C}$ and $\frac{1}{2}(m-2)$ distinct irreducible characters of degree two. So these irreducible characters constitute $4 + 2^2 \left( \frac{1}{2}(m-2) \right) = 4 + 2(m-2) = 2m$ in the sum of squares giving $|D_m|$. It follows that these must be all the irreducible characters.

**Exercises 113.14.**

1. Let $F$ be a subfield of $\mathbb{C}$ and $G$ a finite group such that $F[G]$ is $F$-split. Let $V$ and $W$ be finite dimensional $F[G]$-modules. Show all of the following:
   \[(i) \chi_{V \otimes_F W} = \chi_V \chi_W \quad \text{where} \quad V \otimes_F W \quad \text{is an} \quad F[G]-\text{module with} \quad G-\text{action induced by} \quad g(v \otimes w) = gv \otimes gw, \quad \forall g \in G, \quad v \in V, \quad \text{and} \quad w \in W.\]
   \[(ii) \chi_{V^*} = \overline{\chi}_V.\]
   \[(iii) \chi_{\text{Hom}(V,W)} = \overline{\chi}_V \chi_W.\]

2. Let $F$ be a subfield of $\mathbb{C}$ and $G$ a finite group such that $F[G]$ is $F$-split. Show the characters of $G$ over $F$ form a ring with multiplication given by the previous exercise.

3. Let $F$ be a subfield of $\mathbb{C}$ and $G$ a finite group such that $FG$ is $F$-split. Suppose that $H$ is a subgroup of $G$ and $\chi$ an $F$-character of $H$. Show if $x \in G$, then
   \[\text{ind}^G_H(\chi)(x) = \frac{|Z_G(x)|}{|H|} \sum_{y \in C(x) \cap H} \chi(y)\]
   where $C(x)$ is the conjugacy class of $x$ in $G$.

4. Let $F$ be a subfield of $\mathbb{C}$ and $G$ a finite group such that $F[G]$ is $F$-split. Suppose that $H$ is a subgroup of $G$, $\psi \in \text{class}_F(H)$ and $\psi' \in \text{class}_F(G)$. Show if $\langle \psi', \lambda \rangle = \langle \psi, \text{res}^G_H(\lambda) \rangle$ for all $\lambda \in \text{class}_F(G)$, then $\psi' = \text{ind}^G_H(\psi)$.

5. Let $F$ be a subfield of $\mathbb{C}$ and $G$ a finite group such that $F[G]$ is $F$-split. Suppose that $N$ is a normal subgroup of $G$. Then using the notation of Exercise 110.15(4), show that $1^G_N = \tilde{\chi}_{F[G/N]}$ (where $\chi_{F[G/N]}$ is the regular character of $G/N$).


7. Let $F$ be a field, $G$ a finite group with $H \subset G$ a subgroup. Let $V$ be an $F[G]$-module and $W$ an $F[H]$-module. Then there exists and isomorphism $\text{Hom}_{F[G]}(V, \text{ind}^G_H(W)) \cong \text{Hom}_{F[H]}(\text{res}^G_H(V), W)$.

**114. Torsion Linear Groups**

In this section, we introduce further work of Burnside that focuses on the problem of conditions that force a group $G$ to be finite. Of
course, one would need that the group is a torsion group, i.e., every element has finite order. Burnside looked at the stronger condition on a group $G$, viz., one in which there exist a positive integer $n$ such that $x^n = e$ for every element in the group $G$. Since the group $\times_{i=1}^{\infty} \mathbb{Z}/n\mathbb{Z}$ satisfies this condition yet is an infinite group, other conditions are necessary. Burnside showed if, in addition, such a group $G$ was also a subgroup of $\text{GL}_n(F)$, with $F$ a field of characteristic zero, then $G$ is a finite group. This is the principal result that we prove in this section. This led to a long research project on attempts to generalize it.

We begin with the following lemma:

**Lemma 114.1.** (Frobenius-Schur) Let $F$ be an algebraically closed field and $G$ a (not necessarily finite) group, $\varphi : G \to \text{GL}_n(F)$ an irreducible representation, and $a_{ij} : G \to F$, $1 \leq i, j \leq n$, the coordinate functions of $\varphi$, i.e., $\varphi = (a_{ij})$. Then

1. $\mathbb{M}_n(F) = \langle \varphi(x) \mid x \in G \rangle$.
2. $\mathcal{B} = \{a_{ij} \mid 1 \leq i, j \leq n\}$ is a set of $F$-linearly independent functions.

**Proof.** (1): Extending $\varphi$ to the $F$-algebra map $\varphi : F[G] \to \mathbb{M}_n(F)$ yields (1) by Burnside’s Theorem 107.4.

(2): If $\mathcal{B}$ is $F$-linearly dependent, then it is $F$-linearly dependent as a set of $F$-functionals $F[G] \to F$, say

$$0 = \sum_{i,j} b_{ij} a_{ij}, \quad b_{ij} \in F \text{ not all zero.}$$

Suppose that $b_{i_{0}j_{0}} \neq 0$. Let $e_{i_{0}j_{0}}$ be the $i_{0}j_{0}$th matrix unit. By (1), there exists an $x$ in $F[G]$ such that $\varphi(x) = e_{i_{0}j_{0}}$. Therefore, we have

$$0 = \sum_{i,j} b_{ij} a_{ij}(x) = \sum_{i,j} b_{ij} \delta_{i_{0}i} \delta_{j_{0}j} = b_{i_{0}j_{0}},$$

a contradiction. Hence $\mathcal{B}$ is $F$-linearly independent in $\text{Hom}_F(F[G], F)$.

**Corollary 114.2.** Let $F$ be an algebraically closed field and $G$ a group. Suppose that $\varphi : G \to \text{GL}_n(F)$ is an irreducible representation and $\varphi = (a_{ij})$, $1 \leq i, j \leq n$, the coordinate functions of $G$. Then there exist $x_1, \ldots, x_{n^2}$ elements in $G$ such that $\{x_1, \ldots, x_{n^2}\}$ is an $F$-basis for $\mathbb{M}_n(F)$. Moreover, if $v_k = \varphi(x_l)$, $1 \leq k \leq n^2$, is viewed in $F^{n^2}$, then $\{v_1, \ldots, v_{n^2}\}$ is a basis for $F^{n^2}$.

**Definition 114.3.** Let $G$ be a group. We say that $G$ is of bounded period if there exists a positive integer $N$ such that $x^N = e$ for all elements $x$ in $G$. 
Theorem 114.4. (Burnside) Let $F$ be a field of characteristic zero and $G$ a subgroup of $\text{GL}_n(F)$ (i.e., there exists a faithful representation of $G$ of degree $n$). Suppose that $G$ has bounded period $N$. Then $G$ is a finite group. In fact, $|G| \leq N^{n^3}$.

Proof. We may assume that $F$ is algebraically closed, the inclusion $\iota : G \subset \text{GL}_n(F)$ affords the $F[G]$-module $V = F^n$, and $\iota = (a_{ij})$ with the $a_{ij}$ the coordinate functions of $\iota$. Let $\chi = \chi_V$.

**Case 1.** $V$ is an irreducible $F[G]$-module.

By Corollary 114.2, there exists $\{x_1, \ldots, x_{n^2}\} \subset G$, an $F$-basis for $\mathbb{M}_n(F)$ with $v_k = (a_{ij}(x_k))$, $k = 1, \ldots, n^2$, (viewed in $F^{n^2}$) giving a basis $B = \{v_1, \ldots, v_{n^2}\}$ for $F^{n^2}$. As $\iota$ is a group homomorphism.

\[
(*) \quad \chi(x_k x) = \sum_{i=1}^{n} a_{ii}(x_k x) = \sum_{i=1}^{n} \sum_{j=1}^{n} a_{ij}(x_k) a_{ji}(x), \quad k = 1, \ldots, n^2
\]

are $n^2$-linear equations in the unknowns $a_{ji}(x)$. As every $\chi(x)$, $x \in G$, is a sum of $n$ roots of unity in $\mu_N$, we have

\[
s := |\{\chi(x) \mid x \in G\}| \leq N^n < \infty.
\]

Since $B$ is linearly independent in $F^{n^2}$, the matrix of coefficients of the system $(*$) is invertible, hence by Cramer’s Rule each system $(*$) has a unique solution. It follows that

\[
|G| = |\{a_{ij}(x) \mid x \in G\}| \leq s^{n^2} \leq N^{n^3}.
\]

**Case 2.** $V$ is not an irreducible $F[G]$-module.

As $V$ is reducible, $n > 1$. Since $G$ cannot act transitively on a basis for $V$, there exist representations $\varphi_i : G \to \text{GL}_{n_i}(F)$ with $1 = 1, 2$ each of degree less than $n$, satisfying

\[
\iota(x) = \begin{pmatrix} \varphi_1(x) & \ast \\ 0 & \varphi_2(x) \end{pmatrix} = \begin{pmatrix} \varphi_1(x) & U(x) \\ 0 & \varphi_2(x) \end{pmatrix}
\]

with $U : G \to F^{n_1n_2}$. Set $G_i = \{\varphi_i(x) \mid x \in G\}$, $i = 1, 2$. Then $G_i$ is a group of bounded period $N$ of degree $n_i < n$, so finite with $|G_i| \leq N^{n_i^3}$ for $i = 1, 2$ by induction. Set $H_i = \ker \varphi_i$, then $G_i \cong G/H_i$ and $[G : H_i] < \infty$ for $i = 1, 2$. By Poincaré’s Lemma (Exercise 10.16(6)), we have $[G : H_1 \cap H_2] < \infty$. Let $x \in H_1 \cap H_2$, then

\[
\iota(x) = \begin{pmatrix} I & U(x) \\ 0 & I \end{pmatrix} \quad \text{and} \quad 0 = \iota(x)^N = \begin{pmatrix} I & NU(x) \\ 0 & I \end{pmatrix}.
\]

Consequently, $NU(x) = 0$ and $U(x) = 0$ as $F$ is a field of characteristic zero. Therefore, $H_1 \cap H_2 = 1$, hence $G$ is a finite group and the
representation $\psi : G \to G_1 \times G_2$ given by $x \mapsto (\varphi_1(x), \varphi_2(x))$ is a monomorphism. It follows that

$$|G| \leq |G_1| |G_2| \leq N^{n_1} N^{n_2} \leq N^n.$$  \[ \square \]

**Remark 114.5.** In the above proof we did not needed that $F$ be a field of characteristic zero in Case 1; and in Case 2, we only needed that $\text{char}(F) \nmid N$. In particular, the lemma holds if $V$ is irreducible or if $V$ is reducible and $\text{char}(F) \nmid N$.

Using the proof above, we obtain another theorem of Burnside.

**Theorem 114.6.** (Burnside) Let $F$ be an arbitrary field and $G$ a subgroup of $\text{GL}_n(F)$. Then $G$ is finite if and only if $G$ has finitely many conjugacy classes.

**Proof.** Certainly if $G$ is finite, then it has finitely many conjugacy classes, so we need only show the converse. We use the notation as in the proof of the previous theorem. So $\iota$ is the inclusion map. If the character $\chi$ associated to $\iota$ is irreducible then the proof of Case 1 for Theorem 114 still works as $\{\chi(x) \mid x \in G\}$ is a finite set, which implies that $G$ is a finite group. So we may assume that $\iota$ is not irreducible.

As in the notation of Case 2 in the previous proof, we may assume that $\iota(x) = \left( \begin{array}{cc} \varphi_1(x) & U(x) \\ 0 & \varphi_2(x) \end{array} \right)$, where $G_i = \{\varphi_i(x) \mid x \in G\}$, a representation of degree $n_i < n$, $\varphi : G_i \to \text{GL}_{n_i}(F)$ and $U : G \to F^{n_1 n_2}$, $i = 1, 2$. Let $\psi : G \to G_1 \times G_2$ be given by $x \mapsto (\varphi_1(x), \varphi_2(x))$ as before, and set $H = \ker \psi \cap G$. Then $H \subset \ker \varphi_1 \cap \ker \varphi_2$. Suppose that $\left( \begin{array}{cc} I & U(h) \\ 0 & I \end{array} \right)$ and $\left( \begin{array}{cc} I & U(h') \\ 0 & I \end{array} \right)$ lie in $H$, then

$$\iota(x) = \left( \begin{array}{cc} I & U(h) \\ 0 & I \end{array} \right) \left( \begin{array}{cc} I & U(h') \\ 0 & I \end{array} \right) = \left( \begin{array}{cc} I & U(h + h') \\ 0 & I \end{array} \right).$$

In particular, $H$ is abelian and $H \subset Z_G\left( \begin{array}{cc} I & U(h) \\ 0 & I \end{array} \right)$. It follows that the conjugacy class $C\left( \begin{array}{cc} I & U(h) \\ 0 & I \end{array} \right)$ satisfies

$$|C\left( \begin{array}{cc} I & U(h) \\ 0 & I \end{array} \right)| = \left[ G : Z_G\left( \begin{array}{cc} I & U(h) \\ 0 & I \end{array} \right) \right] \leq |G : H|.$$

As $G$ has finitely many conjugacy classes, so do $G_1$ and $G_2$, hence both are finite by induction. Therefore, $|G : H| \leq |G_1| |G_2| < \infty$. As $G$ has
finitely many conjugacy classes and $H$ is an abelian subgroup of $G$, we have $H$ is also finite. It follows that $G$ is a finite group. □

As mentioned in the introduction to this section, $\times_{i=1}^{\infty} \mathbb{Z}/p\mathbb{Z}$ is an infinite group of bounded period $p$, so certainly not finitely generated. Burnside conjectured in 1902 that any finitely generated torsion group $G$ was finite. This is called the Burnside Conjecture. [Torsion groups are also called periodic groups.] We next improve the Burnside’s theorems in the linear group case.

We need two lemmas.

**Lemma 114.7.** Let $G$ be a finitely generated torsion group. Suppose that $G$ contains an abelian subgroup $H$ of finite index. Then $G$ is a finite group.

**Proof.** Let $G = \bigvee_{i=1}^{m} g_i H$ be a coset decomposition. As $G$ is finitely generated, there exists a finite subset

$$S = \{g_1, \ldots, g_m, g_{m+1}, \ldots, g_n\} \subset G$$

generating $G$ and closed under taking inverses. For each ordered pair $(i, j)$, $1 \leq i, j \leq n$, there exists an integer $r = r(i, j)$, $1 \leq r \leq m$, satisfying

(*) \hspace{1cm} g_i g_j = g_r h_{ij}, \quad h_{ij} \in H.

Set

$$H_0 = \langle h_{ij} \mid h_{ij} \text{ occurring as in (*)} \rangle.$$

Therefore, $H_0$ is a finitely generated group. As $H_0$ is a subgroup of the abelian torsion subgroup $H$, it is finite. Suppose that $(i, j, r)$ is as in (*) and let $1 \leq s \leq n$. Then for each $m = 1, \ldots, m$, there exists a positive integer $v = v(r, s)$ satisfying

$$g_s g_i g_j = g_s g_r h_{ij} = g_v h_{sr} h_{ij} \text{ lies in the coset } g_v H_0.$$

It follows by induction that the group

$$\langle g_1, \ldots, g_n \rangle = \{ \text{words in the } g_i \}$$

lies in $\bigcup_{i=1}^{m} g_i H_0$, hence $G \subset \bigcup_{i=1}^{m} g_i H_0$, so is finite. □

You should compare this proof with the proof of Theorem 16.3. We can use the lemma in the linear group case to see

**Lemma 114.8.** Let $F$ be a field and $G$ a finitely generated torsion subgroup of $\text{GL}_n(F)$. Then $G$ has bounded period. In particular, if $\text{char}(F) = 0$, then $G$ is finite.
Proof. Let $\Delta$ be the prime subfield of $F$. We may assume that $\Delta = \mathbb{Q}$ or $\Delta = \mathbb{Z}/p\mathbb{Z}$, depending on the characteristic of $F$. Since $G$ is finitely generated, we may assume that $F/\Delta$ is a finitely generated field extension. In particular, there exists an intermediate field $F/F_0/\Delta$ with $F_0/\Delta$ a finitely generated purely transcendental field extension and $F/F_0$ a finite field extension. Let $r = [F : F_0]$. The subgroup $G$ of $\text{GL}_n(F)$ acts faithfully on $F^n$, so viewing $F^n$ as $F_0^{nr}$, we have $G$ acts faithfully on the finite dimensional $F_0$-vector space $F_0^{nr}$. In particular, we may assume that $F = F_0$ is a finitely generated purely transcendental extension.

If $x \in G \subset \text{GL}_n(F)$, let $q_x$ be the minimal polynomial of $x$ in $F[t]$. In particular, $\deg q_x \leq \deg f_x = n$, where $f_x$ is the characteristic polynomial of $x$ in $\mathcal{M}_n(F)$.

Case 1. $\Delta = \mathbb{Q}$.

As $x$ in $G$ has finite order, the roots of $q_x$ are roots of unity, so the coefficients of $q_x$ lie in $\mathbb{Z}_{\overline{\mathbb{Q}}} \cap F$, where $\overline{\mathbb{Q}}$ is the algebraic closure of $\mathbb{Q}$ in $\mathbb{C}$ and $\mathbb{Z}_{\overline{\mathbb{Q}}}$ is the integral closure of $\mathbb{Z}$ in $\overline{\mathbb{Q}}$ by Remarks 76.14. Since $F/\mathbb{Q}$ is purely transcendental, $\mathbb{Z}_{\overline{\mathbb{Q}}} \cap F = \mathbb{Z}_{\overline{\mathbb{Q}}} \cap \mathbb{Q} = \mathbb{Z}$ as $\mathbb{Z}$ is an integrally closed domain. It follows that $q_x$ lies in $\mathbb{Z}[t]$ (as each irreducible factor of $q_x$ does) by Proposition 77.4. The coefficients of $q_x$ are elementary symmetric functions in the roots of $q_x$ and these roots are roots of unity, hence the coefficients for each $q_x$ must have bounded absolute value. So there can only be finitely many $q_x$ with roots in a finite set of roots of unity.

Let $\omega$ be a root of $q_x$. Then $m_{\mathbb{Q}}(\omega) \mid q_x \mid f_x$, where $m_{\mathbb{Q}}(\omega)$ is the minimal polynomial of $\omega$ in $\mathbb{C}$. It follows that $\deg m_{\mathbb{Q}}(\omega) \leq n$ by the Cayley-Hamilton Theorem. As a primitive $m$th root of unity $\varepsilon$ in $\mathbb{Q}$ satisfies $\deg m_{\mathbb{Q}}(\varepsilon) = \phi(m)$ and $\phi(m) \to \infty$ as $m \to \infty$, there can only be finitely many $\omega$ that are roots of the $q_x$, all uniformly bounded, i.e., $\{q_x \mid x \in G\}$ is a finite set. Therefore, there exists a positive integer $N$ such that $x^N = 1$ for all $x$ in $G$. In particular, the subgroup $G$ of $\text{GL}_n(F)$ has bounded period. By Burnside’s Theorem 114.4, $G$ is a finite group.

Case 2. $\Delta = \mathbb{Z}/p\mathbb{Z}$ with $p > 0$.

As in Case 1, we have $q_x \in (\mathbb{Z}/p\mathbb{Z})[t]$ satisfies $\deg q_x \leq n$, so there exist finitely many such $q_x$, $x \in G$. The result follows by an analogous argument as in Case 1. □

Theorem 114.9. (Schur) Let $F$ be a field and $G$ a finitely generated torsion subgroup of $\text{GL}_n(F)$. Then $G$ is a finite group.
**Proof.** By Lemma 114.8, \( G \) has bounded period, so by Lemma 114.7, it suffices to show that \( G \) contains an abelian subgroup of finite index. By Remark 114.5, the only part of the proof of Burnside’s Theorem 114.4 in the characteristic zero case that must be modified is the case that the inclusion \( \iota : G \to \text{GL}_n(F) \) is not irreducible. But the proof of Burnside’s Theorem 114.6 produces an abelian subgroup of finite index. \( \square \)

In 1964, Golod-Shafarevich produced an infinite group on three generators in which every element has order a power of a prime, so the Burnside Conjecture is false. A modification of this conjecture was made to the conjecture that any finitely generated torsion group of bounded period is finite. Adrian-Novikov proved this too was false in 1968 (although Novikov claimed to have proven it in 1959). It was modified again to the so-called Restricted Burnside Problem which said: There exist only finitely many \( m \)-generator groups \( G \) of bounded period \( n \). In particular, there exists an integer \( N = N(m, n) \) such that \( |G| < N \) for all such \( G \). This was proven by Zelmanov in 1994 for which he won a Fields Medal.

**Exercise 114.10.** Complete the proof of Case 2 of Theorem 114.9.
CHAPTER XX

Universal Properties and Multilinear Algebra

Given an algebraic object, it is usually very difficult to define a homomorphism from that object to another. Even if we have a well-defined set map, to show that it preserves the structure, i.e., is a homomorphism, we have to show that it preserves all the relations of that object under the map. We saw that free modules were those modules defined by the universal property that a homomorphism from a free module on a basis $\mathcal{B}$ was completely determined by where the basis $\mathcal{B}$ went. This meant that free modules had no non-trivial relations. Many structures in algebra can be defined by such universal properties, we shall call such definitions \textit{categorical definitions}. For example, we also showed that the quotient field of a domain was defined by one (cf. Theorem 26.14). An algebraic object $A$ is defined by a universal property $(\mathcal{P})$ if it satisfies property $(\mathcal{P})$ together with a map (or maps) and if $C$ satisfies property $(\mathcal{P})$, then there is a unique homomorphism(s) $A \to C$ (or $C \to A$) together with a commutative diagram (or commutative diagrams) respecting the given map(s). For example, $F$ is a free $R$-module on a set $\mathcal{B}$ if there is a fixed set map $i : \mathcal{B} \to F$ and if $M$ is an $R$-module and $j : \mathcal{B} \to M$ set map, then there exists a unique $R$-homomorphism $f : F \to M$ satisfying

$$
\begin{array}{ccc}
\mathcal{B} & \xrightarrow{i} & F \\
j & & \downarrow{f} \\
& M &
\end{array}
$$

commutes. This also means that if $F'$ is another free $R$-module on $\mathcal{B}$ with $i' : \mathcal{B} \to F'$ the given map, then there exists a unique $R$-isomorphism $f : F \to F'$. Thus the module $F$ together with the set map $i$ satisfies: there is a canonical bijection

$$
\text{Hom}_{\text{sets}}(\mathcal{B}, N) \to \text{Hom}_R(F, N)
$$

relative to $i : \mathcal{B} \to F$, where the left hand side are set maps and canonical means the the map given uniquely by the commutative diagram.
In this chapter, we shall investigate other universal properties. Many of the details will be left to the reader. As an application we shall define the determinant of an \( R \)-homo- morphism when \( R \) is a commutative ring and satisfying those properties of determinants of matrices that you know.

### 115. Some Universal Properties of Modules

Let \( R \) be a ring. Given an \( R \)-homomorphism \( f : M \to N \) of \( R \)-modules, we know that we should determine its kernel and image. We also saw that cokernels existed. The kernel was a submodule of \( M \) on which \( f \) vanished. Of course, \( f \) can vanish on other submodules. Does \( \ker f \) satisfy a universal property? The answer is yes. Formally we may define the kernel of an \( R \)-homo- morphism as follows:

**Definition 115.1.** Let \( f : A \to B \) be an \( R \)-homomorphism of \( R \)-modules. A kernel of \( f \) is an \( R \)-module \( K \) together with a \( R \)-homomorphism \( i : K \to A \) satisfying

\[
\begin{array}{ccc}
K & \to & 0 \\
\downarrow i & & \downarrow 0 \\
A & \to & B
\end{array}
\]

commutes (where 0 is the zero map), and satisfies the following universal property: If

\[
\begin{array}{ccc}
M & \to & 0 \\
\downarrow g & & \downarrow 0 \\
A & \to & B
\end{array}
\]

is a commutative diagram of \( R \)-modules and \( R \)-homo- morphisms, then there exists a unique \( R \)-homomorphism \( h : M \to K \) satisfying

\[
\begin{array}{ccc}
K & \to & A \\
\downarrow i & & \downarrow f \\
A & \to & B \\
\downarrow h & & \downarrow g \\
M & \to & 0
\end{array}
\]

commutes. Moreover, if the kernel exists, it is unique up to a unique isomorphism.

Our original definition of kernel satisfies the above with \( K = \ker f \) and \( i \) the inclusion map. We have just made a choice of a representation of a kernel object. Since it is unique relative to the inclusion map, there can be no other one. Note also that \( f \) is a monomorphism if and only if the kernel is the zero module (of which there is only one).
We can define a cokernel of an $R$-homomorphism $f : A \to B$ by reversing arrows (called duality), i.e.,

**Definition 115.2.** Let $f : A \to B$ be an $R$-homomorphism of $R$-modules. A cokernel of $f$ is an $R$-module $C$ together with a $R$-homomorphism $j : B \to C$ satisfying

$$
\begin{array}{c}
A \xrightarrow{f} B \\
\downarrow{0} \quad \downarrow{j} \\
C
\end{array}
$$

commutes, and satisfying the following universal property: If

$$
\begin{array}{c}
A \xrightarrow{f} B \\
\downarrow{0} \quad \downarrow{g} \\
N
\end{array}
$$

is a commutative diagram of $R$-modules and $R$-homomorphisms, then there exists a unique $R$-homomorphism $h : C \to N$ satisfying

$$
\begin{array}{c}
A \xrightarrow{f} B \xrightarrow{j} C \\
\downarrow{0} \quad \downarrow{g} \quad \downarrow{h} \\
N
\end{array}
$$

commutes. Moreover, if the cokernel exists, it is unique up to a unique isomorphism.

Then our definition of cokernel is the unique quotient object $B/\text{im } f$ relative to the canonical epimorphism. One problem here is that we want to define the image by a universal property, and it will depend on the definition of cokernel. Note that it is automatic that the map to a cokernel is surjective, so isomorphic to a quotient of $B$ by some submodule. We define the image of an $R$-homomorphism as follows:

**Definition 115.3.** Let $f : A \to B$ be an $R$-homomorphism of $R$-modules. An image of $f$ is kernel of $j : B \to C$, where $C$ is a cokernel.

**Remarks 115.4.** By choice of nomenclature, a cokernel of a kernel of an $R$-homomorphism $f : A \to B$, should be called a coimage of $f : A \to B$. Therefore, $A/\ker f$ is the coimage of $f$. By the First Isomorphism Theorem, the quotient $A/\ker f$ is (canonically) isomorphic to the image of $f$, i.e., we can view the First Isomorphism Theorem as saying the coimage of $f$ and the image of $f$ are canonically isomorphic. These notions are important in category theory when they hold.
When we gave examples of modules, two that arose were the direct product and direct sum of modules. We now define these by categorical definitions.

**Definition 115.5.** Let \( A \) and \( B \) be \( R \)-modules. Then a *product* of \( A \) and \( B \) is an \( R \)-module \( P \) together with \( R \)-homomorphisms

\[
P \xrightarrow{\pi_A} A \\
\pi_B \\
B
\]

satisfying the following universal property: If

\[
X \xrightarrow{f} A \\
g \\
\downarrow \\
\downarrow \\
B
\]

are \( R \)-homomorphisms, then there exists a unique \( R \)-homomorphism \( h : X \to P \) such that

\[
\begin{array}{c}
X \\
\downarrow \\
P \\
\downarrow \\
A \\
\downarrow \\
\pi_B \\
\downarrow \\
B
\end{array}
\]

commutes. Moreover, if the product exists, it is unique up to a unique isomorphism and denoted by \( A \prod B \).

Of course, this product is the external direct product previously defined via the projection maps.

The *coproduct* \( Y \) of two \( R \)-modules \( A \) and \( B \) is the dual notion of product defined by reversing arrows, i.e., it is a diagram

\[
A \xrightarrow{\iota_A} Y \\
\downarrow \\
\downarrow \\
\iota_B \\
B
\]

satisfying the obvious universal property. It is unique up to a unique isomorphism and denoted by \( A \coprod B \). Of course, this coproduct is the external direct sum previously defined via the inclusion maps of coordinates.
One defines a product and a coproduct of an arbitrary collection \( \{ M_i \}_{i \in I} \) of \( R \)-modules in the obvious way and denotes it by \( \prod_I M_i \) and \( \coprod_I M_i \) respectively. For all \( j \), they come equipped with both epimorphisms \( \pi_j : \prod_I M_i \rightarrow M_j \) and monomorphisms \( \iota_j : M_j \rightarrow \coprod_I M_i \) respectively which we identify as the ones that previously arose. If \( I \) is finite these are isomorphic and are identified as before.

As an example of how to use universal properties, we show:

**Lemma 115.6.** An \( R \)-module \( A \) is a coproduct of the \( R \)-submodules \( A_1, \ldots, A_n \) if and only if there exist \( R \)-homomorphisms \( \iota_j : A_j \rightarrow A \) and \( \pi_j : A \rightarrow A_j \) for \( j, k = 1, \ldots, n \) satisfying \( \pi_k \iota_j = \delta_{kj} 1_{A_j} \) and \( \sum_j \iota_j \pi_j = 1_A \).

**Proof.** (\( \Rightarrow \)): The definition of coproduct gives rise to the \( \iota_j \)'s. For a fixed \( k \) and and each \( j \), we have a diagram

\[
\begin{array}{ccc}
A_j & \xrightarrow{\delta_{kj} 1_{A_j}} & A_k \\
\downarrow{\iota_j} & & \downarrow{A} \\
A & & \end{array}
\]

By the universal property of coproduct, there exists a unique \( R \)-homomorphism \( \pi_k : A \rightarrow A_k \) such that

\[
\begin{array}{ccc}
A_j & \xrightarrow{\delta_{kj} 1_{A_j}} & A_k \\
\downarrow{\iota_j} & & \downarrow{\pi_k} \\
A & & \end{array}
\]

commutes. By distributivity, \( (\iota_1 \pi_1 + \cdots + \iota_n \pi_n) \iota_j = \iota_j \). Then

\[
\begin{array}{ccc}
A_j & \xrightarrow{\iota_j} & A \\
\downarrow{\iota_j} & & \downarrow{1_A} \\
A & & \end{array} \quad \text{and} \quad \begin{array}{ccc}
A_j & \xrightarrow{\iota_j} & A \\
\downarrow{\iota_j} & & \downarrow{\iota_1 \pi_1 + \cdots + \iota_n \pi_n} \\
A & & \end{array}
\]

both commute for each \( j \). Uniqueness gives \( 1_A = \iota_1 \pi_1 + \cdots + \iota_n \pi_n \).

(\( \Leftarrow \)): If \( f_j : A_j \rightarrow B \) is an \( R \)-homomorphisms for each \( j \), then defining \( f : A \rightarrow B \) by \( f = \sum_i f_i \pi_i \) works. \( \square \)

Suppose that we have \( R \)-homomorphisms of \( R \)-modules

\[
A \xrightarrow{f} B \xrightarrow{h} D \quad \text{and} \quad A \xrightarrow{g} C \xrightarrow{j} D,
\]
then by the universal properties of product and coproduct, there exist unique $R$-homomorphisms

$$A \xrightarrow{(f,g)} B \coprod C \xrightarrow{h \coprod j} D$$

given by

$$(f, g) := \iota_B f + \iota_C g$$

$$h \coprod j := h\pi_B + j\pi_C,$$

i.e., as we are identifying $A \prod B$ and $A \coprod B$, i.e., we have a diagram

In particular, we have

$$(h \coprod g)(f, g) = hf + jg.$$  

An interesting example of this is when $A = B = C$ and $f = g = 1_A$. Set $\Delta = (1_A, 1_A) : A \to A \coprod A$, called the diagonal map. Then the above shows

$$A \xrightarrow{h + j} D$$

$$A \coprod A \xrightarrow{h \coprod j} D.$$  

commutes. Note that this commutative diagram categorically determines addition in $\text{Hom}_R(A, D)$.

**Exercises 115.7.**

1. Define a free group, i.e., a group satisfying the obvious definition. Do you know of any examples of such an item.
2. Define the product and coproduct of a collection of $R$-modules $\{M_i \mid i \in I\}$ and show that these are the external direct product and external direct sum of $R$-modules.
3. Define the product and coproduct of two groups. Can you identify them?

4. Let \( f : M \to N \) be an \( R \)-homomorphism of \( R \)-modules. Show that there exists a natural \( R \)-homomorphism \( \text{coim} f \to \text{im} f \) arising from universal properties without using the First Isomorphism Theorem. Then show that it is an \( R \)-isomorphism.

5. Let \( R \) be a commutative ring, \( S \subset R \) a multiplicative set, and \( M \) an \( R \)-module. Let \( S^{-1}M := \{ \frac{m}{s} \mid m \in M, s \in S \} \). Define an addition and \( R \)-action on \( S^{-1}M \) making it an \( R \)-module and an \( S^{-1}R \)-module. This is called the localization of \( M \) at \( S \). Determine the universal property that such a module satisfies.

116. Tensor Products

Let \( V \) be a vector space over \( F \). If \( v \) and \( w \) are vectors in \( V \), in general, there is no way to define a product of \( v \) and \( w \) in \( V \). We want to rectify this, by creating a vector space in which a ‘product’ of \( v \) and \( w \) makes sense, called the tensor product. We shall do this in generality based upon a universal property. For vector spaces, as they are free modules, what we get is what one desires, but in general, unexpected results can occur.

Let \( M_1, \ldots, M_n \), and \( N \) be \( R \)-modules and \( M_1 \times \cdots \times M_n \) the cartesian product of the \( M_i \). A map

\[
f : M_1 \times \cdots \times M_n \to N
\]

is called (\( R \)-) \( n \)-linear or (\( R \)-) \( n \)-multilinear if \( f \) is linear in each variable, i.e.,

\[
f(x_1, \ldots, rx_i+x_i', \ldots, x_n) = rf(x_1, \ldots, x_i, \ldots, x_n) + f(x_1, \ldots, x_i', \ldots, x_n)
\]

for all \( x_j \in M_j \), \( x_i' \in M_i \), and \( r \in R \). If \( n = 2 \) an \( n \)-linear form is called an \( R \)-bilinear form. If \( R \) is a commutative ring we would like to convert bilinear maps into \( R \)-homomorphisms. This is done as follows:

**Definition 116.1.** Let \( R \) be a commutative ring and \( M, N \) two \( R \)-modules. An \( R \)-module \( T \) is called a tensor product of \( M \) and \( N \) if there exists an \( R \)-bilinear map \( \iota : M \times N \to T \) satisfying the following universal property: If \( A \) is an \( R \)-module and \( j : M \times N \to A \) an \( R \)-bilinear map, then there exists a unique \( R \)-homomorphism \( f : T \to A \) such that

\[
\begin{array}{ccc}
M \times N & \xrightarrow{\iota} & T \\
\downarrow j & & \downarrow f \\
A & & A
\end{array}
\]
commutes. If a tensor product \( \iota : M \times N \to T \) exists, it is unique up to a unique isomorphism and \( T \) is denoted by \( M \otimes_R N \).

**Lemma 116.2.** Let \( R \) be a commutative ring and \( M, N \) two \( R \)-modules. Then a tensor product of \( M \) and \( N \) exists.

**Proof.** Let \( P \) be the free \( R \)-module on basis \( \mathcal{B} = M \times N \). Let \( W \) be the submodule of \( P \) generated by the following:

\[
\begin{align*}
(m + m', n) - (m, n) - (m', n) \\
(m, n + n') - (m, n) - (m, n') \\
(rm, n) - r(m, n) \\
(m, rn) - r(m, n)
\end{align*}
\]

for all \( m, m' \) in \( M \), \( n, n' \) in \( N \), and \( r \) in \( R \). Let \( \varpi : P \to P/W \) be the canonical \( R \)-epimorphism. Then the composition \( \iota \) of the maps \( M \times N \xrightarrow{\text{inc}} P \xrightarrow{\varpi} P/W \) is \( R \)-bilinear. Set \( T = P/W \) and \( M \times N \to T \). We denote \( \iota(m, n) \) by \( m \otimes n \). Then every element in \( T \) is a finite sum

\[
\sum_{\text{finite}} r_i (m_i \otimes n_i) = \sum_{\text{finite}} (r_i m_i) \otimes n_i = \sum_{\text{finite}} m_i \otimes (r_i n_i),
\]

for appropriate elements \( m_i \) in \( M \), \( n_i \) in \( N \), and \( r_i \) in \( R \), as \( r(m \otimes n) = (rm) \otimes n = m \otimes (rn) \) for all \( m \) in \( M \), \( n \) in \( N \), and \( r \) in \( R \). Suppose that \( j : M \times N \to A \) is \( R \)-bilinear. Define \( f : T \to A \) as follows: By the universal property of freeness on basis \( M \times N \), there exists a unique \( R \)-homomorphism \( g : P \to A \) satisfying

\[
\begin{array}{ccc}
M \times N & \xrightarrow{\iota} & P \\
\downarrow j & & \downarrow g \\
M \times N & \xrightarrow{j} & A
\end{array}
\]

commutes. Now, with \( \varpi : P \to P/W \) the canonical \( R \)-epimorphism,

\[
\begin{array}{ccc}
M \times N & \xrightarrow{j} & P \xrightarrow{\varpi} P/W \\
\downarrow j & & \downarrow g \\
M \times N & \xrightarrow{j} & A
\end{array}
\]

and clearly, \( g|_{\ker \varpi} = 0 \), i.e., \( g|_W = 0 \) as \( j \) is \( R \)-bilinear. Hence

\[
\begin{array}{ccc}
0 & \xrightarrow{\text{inc}} & P \xrightarrow{-} P/W \\
\downarrow 0 & & \downarrow g \\
0 & \xrightarrow{-} & A
\end{array}
\]
is exact. But \( P/W \) is the cokernel of the inclusion map \( \text{inc} : W \to P \), so there exists a unique \( R \)-homomorphism \( \overline{g} : P/W \to A \) satisfying

\[
\begin{array}{c}
P \xrightarrow{g} P/W \\
\downarrow \quad \quad \downarrow \\
A \quad \quad \quad \overline{g}
\end{array}
\]

commutes. As

\[
\begin{array}{c}
M \times N \xrightarrow{\iota} P/W \\
\downarrow j \quad \quad \downarrow \overline{g} \\
A
\end{array}
\]

commutes and \( \overline{g} \) is unique, we have established existence. [If you do not like using the universal property of cokernels, you can prove this directly, i.e., if \( A \) is a submodule of an \( R \)-module \( B \) and \( f : B \to C \) is an \( R \)-homomorphism satisfying \( A \subset \ker \varphi \), then there exists a unique \( R \)-homomorphism \( \overline{\varphi} : B/A \to C \) satisfying

\[
\begin{array}{c}
B \xrightarrow{\varphi} C. \\
\downarrow \quad \quad \downarrow \overline{\varphi} \\
B/A
\end{array}
\]

commutes.] We leave the proof of uniqueness to the reader. \( \Box \)

We adopt the notation above, i.e., elements of \( M \otimes_R N \) are finite sums \( \sum r_i (m_i \otimes n_i) \), etc.

Remark 116.3. Let \( R \) be a commutative ring.

1. If \( M, N \), and \( A \) are \( R \)-modules, by the universal property of tensor products, the module structure of \( M \otimes_R N \to A \) is determined by a \( R \)-bilinear map \( M \times N \to A \) (hence on generators \( (m, n) \) in \( M \times N \)). In particular, an \( R \)-homomorphism \( \varphi : M \otimes_R N \to A \), is completely determined by the \( \varphi(m \otimes n) \).

2. Let \( M_i, N_i \) be \( R \)-modules, and \( \varphi_i : M_i \to N_i \) an \( R \)-homomorphism for \( i = 1, 2 \). By the universal property of tensor products, there exists a unique \( R \)-homomorphism \( \varphi_1 \otimes \varphi_2 : M_1 \otimes_R M_2 \to N_1 \otimes_R N_2 \) induced by \( m_1 \otimes m_2 \mapsto \varphi(m_1) \otimes \varphi(m_2) \) as \( M_1 \times M_2 \to N_1 \otimes_R N_2 \) given by \( (m_1, m_2) \mapsto \varphi(m_1) \otimes \varphi(m_2) \) is \( R \)-bilinear.
If $R$ is a commutative ring, then the lemma shows that there exists a bijection of sets

$$\{ f : M \times N \to A \mid f \text{ an } R\text{-bilinear map} \}$$

$$\to \{ f : M \otimes_R N \to A \mid f \text{ an } R\text{-homomorphism} \}$$

i.e., $\otimes_R$ converts $R$-bilinear maps to $R$-linear maps, i.e., $R$-homomorphisms, as desired.

**Properties 116.4.** Let $R$ be a commutative ring and $P, M, N$ $R$-modules.

1. $0 \otimes 0 = m \otimes 0 = 0 \otimes n$ is the zero of $M \otimes_R N$ for any $m$ in $M$, any $n$ in $N$.

2. The map $\iota_M : R \otimes_R M \to M$ induced by $r \otimes m \mapsto rm$ is a canonical $R$-homomorphism, where canonical means if $f : M \to N$ is an $R$-homomorphism, then we have commutative diagram

$$\begin{array}{ccc}
R \otimes_R M & \xrightarrow{\iota_M} & M \\
\downarrow^{1_R \otimes \varphi} & & \downarrow^\varphi \\
R \otimes_R N & \xrightarrow{\iota_N} & N.
\end{array}$$

[Why is it well defined?] Check that $M \to R \otimes_R M$ given by $m \mapsto 1 \otimes m$ is the inverse.

3. Let $M_i, i \in I, N$ be $R$-homomorphisms. Then there exists a canonical $R$-iso-morphism

$$\prod_I (M_i \otimes_R N) \to (\prod_I M_i) \otimes_R N,$$

i.e., $\prod_I$ and $\otimes_R$ commute. (Similarly, on the other side.):

Since the map

$$\prod_I M_i \times N \to \prod_I (M_i \otimes_R N_i)$$
given by $\{\{m_i\}_I, n\} \mapsto \{m_i \otimes n\}_I$

is $R$-bilinear, the universal property of the tensor products induces a unique $R$-homomorphism $f : \left(\prod_I M_i\right) \otimes_R N \to \prod_I (M_i \otimes_R N_i)$. For each $j \in I$, we have an $R$-bilinear map $M_j \times N \to \left(\prod_I M_i\right) \otimes_R N$ given by $(m, n) \mapsto \{\delta_{ij} m_i\} \otimes n$, inducing a unique $R$-homomorphism $M_j \otimes_R N \to \left(\prod_I M_i\right) \otimes_R N$ by the universal property of tensor products. Hence, by the universal property of coproducts, there exists a unique $R$-homomorphism $g : \left(\prod_I M_i\right) \otimes_R N \to \prod_I M_i \otimes_R N$. Check that $f$ and $g$ are inverses of each other.
4. $M \otimes_R N$ is canonically isomorphic to $N \otimes_R M$ induced by the $R$-bilinear map $M \times N \rightarrow N \otimes_R M$ given by $(m, n) \mapsto n \otimes m$.

5. $(M \otimes_R N) \otimes_R P$ is canonically isomorphic to $M \otimes_R (N \otimes_R P)$.

6. $(\bigsqcup_i R) \otimes_R M \cong \bigsqcup_i (R \otimes_R M) \cong \bigsqcup_i M$ with all $R$-isomorphisms canonical.

7. Suppose that $M$ is a free $R$-module on basis $\mathcal{B}$ and $N$ is a free $R$-module on basis $\mathcal{C}$. Then $M \otimes_R N$ is a free $R$-module on basis $\{b \otimes c \mid b \in \mathcal{B}, c \in \mathcal{C}\}$. In particular,

$$\text{rank}_R(M \otimes_R N) = \text{rank}_R M \cdot \text{rank}_R N :$$

As $R \cong Rx$ for all $x$ in $\mathcal{B}$ or $\mathcal{C}$, we have $M \cong \bigsqcup_\mathcal{B} R$ and $N \cong \bigsqcup_\mathcal{C} R$. Therefore,

$$M \otimes_R N \cong (\bigsqcup_\mathcal{B} R) \otimes_R (\bigsqcup_\mathcal{C} R) \cong \bigsqcup_\mathcal{B} (R \otimes_R (\bigsqcup_\mathcal{C} R)) \cong \bigsqcup_\mathcal{B} \bigsqcup_\mathcal{C} R.$$

8. If $M$, $N$ are not free $R$-modules, bad things can occur. For example, the abelian group $\mathbb{Z}/2\mathbb{Z} \otimes_{\mathbb{Z}} \mathbb{Z}/3\mathbb{Z}$ is not free. For if $[a]_n$ is the congruence class of the integer $a$ modulo $n$, then $[1]_2 \otimes [1]_3 = [1]_2 \otimes 4[1]_3 = 4[1]_2 \otimes [1]_3 = 0$. Similarly, $[1]_2 \otimes [-1]_3 = 0$.

9. (Base Extension) Let $\varphi : R \rightarrow T$ be a ring homomorphism of commutative rings. We know that any $T$-module $B$ becomes an $R$-module via the pullback, i.e., $rx := \varphi(r)x$ for all $r \in R$, $x \in B$, i.e., via $\varphi$. In particular, $T$ becomes an $R$-module, hence so does $T \otimes_R M$. We can now make $T \otimes_R M$ into a $T$-module by defining

$$\alpha(\beta \otimes m) := (\alpha \beta) \otimes m,$$

for all $\alpha$ and $\beta$ in $T$ and $m$ in $M$. We call the $T$-module $T \otimes_R M$ the base extension of the $R$-module $M$ to a $T$-module (via $\varphi$).

**Examples 116.5.** Let $K/F$ be a field extension. Then $K$ is an $F$-vector space, say on basis $\mathcal{C}$. Let $V$ be an $F$-vector space on basis $\mathcal{B}$. The $K \otimes_F V$ is an $F$-vector space on basis $\{c \otimes b \mid c \in \mathcal{C}, b \in \mathcal{B}\}$, so $\dim_F K \otimes_F V = \dim_F K \dim_F V$ and $K \otimes_F V$ is a $K$-vector space on basis $\{1 \otimes b \mid b \in \mathcal{B}\}$. For example, $\mathbb{C} \otimes_{\mathbb{R}} \mathbb{C}$ is a four dimensional real vector space and a 2-dimensional complex vector space.

**Question 116.6.** Can you make $\mathbb{C} \otimes_{\mathbb{R}} \mathbb{C}$ into a ring? If so, is it a field?, a domain?, commutative?

**Proposition 116.7.** Let $V$ and $W$ be finite dimensional $F$-vector spaces, $V^* = \text{Hom}_F(V, F)$, the dual space of $V$ (or finitely generated $R$-free modules with $R$ commutative). Then the natural map

$$\varphi : V^* \otimes_F W \rightarrow \text{Hom}_F(V, W)$$
induced by \( f \otimes w \mapsto (v \mapsto f(v)w) \) is an \( F \)-linear isomorphism.

**Proof.** Although the map in the proposition is independent of any basis, we use bases to prove it. Let \( \mathcal{B} = \{v_1, \ldots, v_n\} \) be an \( F \)-basis for \( V \) and \( \mathcal{B}^* = \{f_1, \ldots, f_n\} \) the dual \( F \)-basis for \( V^* \), i.e., \( f_i(v_j) = \delta_{ij} \), so \( f_i \) is the coordinate function on \( v_i \). Therefore, if \( v \) is an element of \( V \), \( v = \sum_i f_i(v) v_i \). Define

\[
\psi : \text{Hom}_F(V, W) \to V^* \otimes_F W \text{ by } T \mapsto \sum_i f_i \otimes T v_i.
\]

**Check:** \( \varphi \) and \( \psi \) are both \( F \)-linear. Then we have

\[
\varphi \psi T(v) = \varphi \left( \sum f_i \otimes T v_i \right)(v) = \sum f_i(v) T v_i = \sum T \left( f_i(v) v_i \right) = T(v),
\]

so \( \varphi \psi = 1_{\text{Hom}_F(V, W)} \). Hence \( \varphi \) is an epimorphism. As \( \dim_F (V^* \otimes_F W) = \dim_F \text{Hom}_F(V, W) < \infty \), \( \varphi \) is an \( F \)-isomorphism (as both are finite dimensional \( F \)-vector spaces).

[If \( V \) and \( W \) are only assumed to be finitely generated free \( R \)-modules, one must show that \( \psi \varphi = 1_{V^* \otimes_R W} \).] \( \square \)

**Remarks 116.8.** We look at an interesting special case of the proposition above. Let \( V = W \) and \( \mathcal{B}, \mathcal{B}^* \) as in the proposition. Then the proposition gives an \( F \)-linear transformation \( \varphi : V^* \otimes_F V \to \text{End}_F(V) \).

Moreover, the inverse map \( \psi \) satisfies \( \psi(1_{\text{End}_F(V)}) = \sum f_i \otimes v_i \). Let

\[ \text{End}_F(V) \xrightarrow{\psi} V^* \otimes_F V \xrightarrow{\rho} F, \]

where \( \rho \) is induced by \( f \otimes v \mapsto f(v) \). It follows that for \( T \in \text{End}_F(V) \), we have

\[ \rho \psi(T) = \rho \left( \sum f_i \otimes T v_i \right) = \sum f_i (T v_i). \]

In particular, if \( T v_i = \sum \alpha_{ji} v_j \), then

\[ \rho \psi(T) = \sum_i f_i \left( \sum_j \alpha_{ji} v_j \right) = \sum_i \alpha_{ji} f_i(v_j) = \sum_i \alpha_{ii} = \text{trace } [T]_{\mathcal{B}}. \]

Since \( \varphi \) is defined independently of \( \mathcal{B} \) and \( \psi \) is its inverse, \( \psi \) is also independent of \( \mathcal{B} \); therefore, \( \rho \psi(T) \) is independent of \( \mathcal{B} \), i.e., the trace of \( T \) defined by

\[ \text{trace } T := \rho \psi(T) = \text{trace } [T]_{\mathcal{B}} \]
is independent of \( \mathcal{B} \). Since \( \psi(1_{\text{End}_F(V)}) = \sum f_i \otimes v_i \), we have \( \varphi(f_i \otimes v_i) = 1_{\text{End}_F(V)} \). The element \( \sum f_i \otimes v_i \) is called the \textit{Casimir element}. It too is independent of \( \mathcal{B} \) and \( \mathcal{B}^* \).

It is useful to generalize the tensor product of modules over a commutative ring to the tensor product of modules over an arbitrary ring, i.e., not necessarily commutative. We indicate this, leaving the proofs as exercises. We first must modify the definition of a bilinear map.

\begin{definition}
Let \( R \) be a ring with \( M \) a right \( R \)-module, \( N \) a left \( R \)-module, and \( A \) an abelian group. A map \( j : M \times N \to A \) is called a \textit{\( R \)-balanced biadditive map}, if it is additive in each variable and satisfies \( j(mr, n) = j(m, rn) \) for all \( r \in R, m \in M, \) and \( n \in N \).
\end{definition}

\begin{definition}
Let \( R \) be a ring with \( M \) a right \( R \)-module, \( N \) a left \( R \)-module. An abelian group \( T \) is then called a \textit{tensor product} of \( M \) and \( N \) if there exists an \( R \)-balanced biadditive group homomorphism \( \iota : M \times N \to T \), i.e., \( \iota(mr, n) = \iota(m, rn) \), for all \( m \in M, n \in N, r \in R \), satisfying the following universal property: If \( A \) is an abelian group and \( j : M \times N \to A \) an \( R \)-balanced biadditive map, then there exists a unique balanced group homomorphism \( f : T \to A \) such that

\[
\begin{array}{ccc}
M \times N & \xrightarrow{j} & T \\
\downarrow & & \downarrow f \\
A & & A
\end{array}
\]

commutes.

\begin{theorem}
Let \( R \) be a ring with \( M \) a right \( R \)-module, \( N \) a left \( R \)-module. Then the tensor product of \( \iota : M \times T \) exists and is unique up to a unique isomorphism. We denote it by \( M \otimes_R N \).
\end{theorem}

\begin{remark}
Let \( Z \) be the center of \( R \) and \( \mathfrak{A} \) the ideal in \( M \otimes_Z N \) generated by \( \{mr \otimes n - m \otimes rn \mid m \in M, n \in N, r \in R \} \). Then \( M \otimes_R N \cong (M \otimes_Z N)/\mathfrak{A} \).
\end{remark}

\begin{definition}
Let \( R \) and \( S \) be rings. An abelian group \( M \) is called an \textit{(\( R - S \))-bimodule} if \( M \) is a left \( R \)-module and a right \( S \)-module satisfying \( r(ms) = (rm)s \) for all \( s \in S, r \in R, \) and \( m \in M \). We write \( M =_R M_S \) if \( M \) is an \( (R - M) \)-bimodule. It is also convenient to write \( sN \) if \( N \) is a left \( S \)-module and \( N_S \) if \( N_S \) is a right \( S \)-module.

For example, if \( S \) is an \( R \)-algebra (so \( R \) is commutative), then, by definition, \( S \) is an \( (R - R) \)-bimodule (as well as an \( (S - S) \)-bimodule).
\end{definition}

\begin{remarks}
Let \( R, S \) be rings, \( M \) a right \( S \)-module and \( N \) a left \( S \)-module.
\end{remarks}
1. Suppose that \( M = {}_RM_S \) is an \((R-S)\)-bimodule, then the abelian group \( M \otimes_S N \) is a left \( R \)-module via \( r(r' \otimes x) = rr' \otimes x \) for all \( r, r' \in S, x \in N \).

2. (Base Extension) Let \( \varphi : S \to R \) be a ring homomorphism. Then the map \( \psi : R \otimes_R S \to R \) induced by \( r \otimes s \mapsto rs(= r\varphi(s)) \) gives \( R \) the structure of an \((R-S)\)-bimodule. In particular, the map \( \varphi \) makes \( R \otimes_S N \) into a left \( R \)-module.

3. Let \( \varphi : S \to R \) be a ring homomorphism. Then \( R \otimes_S S \cong R \) as left \( R \)-modules as follows: Let \( \psi : R \times S \to R \) be the right \( S \)-homomorphism defined by \( (r, s) \mapsto rs(= r\varphi(s)) \). Then \( \psi \) is easily checked to be a balanced map, e.g., \( \psi(rs, s') = (rs)s' = r(ss') = \psi(r, ss') \) for all \( r \in R, s, s' \in S \). In particular, it induces a group homomorphism \( \Psi : R \otimes_S S \to R \) satisfying \( \Psi(r \otimes s) = rs \). The map \( \Psi \) is a left \( R \)-homomorphism, since \( \Psi(r'(r \otimes s)) = \Psi(r'r \otimes s) = r'rs = r\Psi(r \otimes s) \) for all \( r', r \in R, s \in S \). The \( R \)-homomorphism \( \Phi : R \to R \otimes_S S \) induced by \( r \mapsto r \otimes 1 \) is the inverse of \( \Psi \), since \( \Psi \circ \Phi(r) = \Psi(r \otimes 1) = r \) and \( \Phi \circ \Psi(r \otimes s) = \Phi(rs) = rs \otimes 1 = r \otimes s \).

This means that \( \Phi \) is the inverse on the generators of \( R \otimes_S S \), hence on all of \( R \otimes_S S \).

4. If \( M = {}_sM_S \) is an \((S-S)\)-bimodule then \( R \otimes_S (M \otimes_S N) \) is an \( R \)-module by \( r(r' \otimes m \otimes n) = rr' \otimes m \otimes n \) for all \( r, r' \in R, m \in M, n \in N \).

5. If \( M_i, i \in I \), are right \( S \)-modules and \( N_i, i \in I \), are left \( S \)-modules, then \( \prod_I (A_i \otimes_S B_i) \cong (\prod_I A_i) \otimes_S \prod_I B_i \), as abelian groups (respectively, as left \( R \)-modules if \( M_i = {}_RM_i \) or as \( R \)-modules if \( {}_SN_i \)) for all \( i \in I \), i.e., and \( \otimes \) commute. In particular, if \( N \) is \( S \)-free, i.e., \( N \cong S^I \) as a left \( S \)-free module, then \( R \otimes_S N \cong R^I \).

We also need the structure on Hom-sets.

**Remark 116.15.** Let \( R \) and \( S \) be rings. If \( M \) and \( N \) are left \( R \)-modules then \( \text{Hom}_R(M, N) \) is an abelian group. (It is usually not an \( R \)-homomorphism since \( rr' \neq rr' \). If, however, \( M = {}_RA_S \) is an \((R-S)\)-bimodule, then \( \text{Hom}_R(M, N) \) is a left \( S \)-module. Indeed for each \( s \in S \), and \( \varphi \in \text{Hom}_R(M, N) \), define \( {}_s\varphi((m) = \varphi(ms) \) for all \( x \in M, s \in S \). Then \( {}_s\varphi \) is a left \( R \)-homomorphism, hence induces a left \( R \)-action on \( \text{Hom}_R(M, N) \). This \( S \)-action makes \( \text{Hom}_R(M, N) \) into a left \( S \)-module.

**Theorem 116.16.** (Adjoint Associativity Theorem) Let \( R \) and \( S \) be rings, \( A = {}_RA_R \) a right \( R \)-module, \( B = {}_RB_S \) an \((R-S)\)-bimodule, and \( C = {}_sC \) a right \( S \)-module. Then there is an isomorphism of abelian groups,

\[
\text{Hom}_S((A \otimes_R B, C), C) \cong \text{Hom}_R(A, \text{Hom}_S(A, (B, C))).
\]

We say that \( A \otimes_R - \) and \( \text{Hom}_S(A, -) \) are adjoints.
**Proof.** Let $\varphi : A \otimes_R B \to C$ be an $S$-homomorphism. For each $a \in A$, define a map $\Phi(a) : B \to C$ by $\Phi(a)(b) = \varphi(a \otimes b)$. Check that $\Phi(a)$ is $S$-homomorphism of right $S$-modules and the map $\Phi : A \to \Hom_S(B, C)$ given by $a \mapsto \varphi(a \otimes b)$ is an $R$-homomorphism of right $R$-modules. Define $f : \Hom_S(A \otimes_R, BC) \to \Hom_R(A, \Hom_S(B, C))$ by $f(\varphi) = \Phi$. Now let $\Psi : A \to \Hom_S(B, C)$ be a group homomorphism. The map $\psi : A \times B \to C$ defined by $(a, c) \mapsto \Psi(a)(c)$ is an $R$-balanced map, so induces a group homomorphism $\psi : A \otimes_R B \to C$. Then $g : \Hom_R(A, \Hom_S(A, (B, C))) \to \Hom_S((A \otimes_R B, C), C)$ defined by $g(\Psi) = \psi$ is a group homomorphism and the inverse of $f$. □

**Exercises 116.17.**

1. Prove that the tensor product of two modules over a commutative ring is unique up to a unique isomorphism.

2. Show $0 \otimes 0 = m \otimes 0 = 0 \otimes n$ is the zero of $M \otimes_R N$ for any $m$ in $M$, any $n$ in $N$.


4. Let $m$ and $n$ be positive integer and $d$ the greatest common divisor of $m$ and $n$. Then $\mathbb{Z}/m\mathbb{Z} \otimes_{\mathbb{Z}} \mathbb{Z}/n\mathbb{Z} \cong \mathbb{Z}/d\mathbb{Z}$ (but not canonically).

5. In the proof of Proposition 116.7, show that $\psi \varphi = 1_{V^* \otimes_R W}$ if $V$ and $W$ are finitely generated free $R$-modules.

6. Let $R$ be a commutative ring, $M$ an $R$-module, and

$$0 \to A \xrightarrow{f} B \xrightarrow{g} C \to 0$$

a short exact sequence of $R$-modules. Show that

$$A \otimes_R M \xrightarrow{f \otimes 1_M} B \otimes_R M \xrightarrow{g \otimes 1_M} C \otimes_R M \to 0$$

and

$$M \otimes_R A \xrightarrow{1_M \otimes f} M \otimes_R B \xrightarrow{1_M \otimes g} M \otimes_R C \to 0$$

are exact.

7. Let $R$ be a commutative ring, $S \subset R$ a multiplicative set, and $M$ an $R$-module. Show that there exists a natural $S^{-1}R$-isomorphism $S^{-1}M \to S^{-1}R \otimes_R M$. (Cf. Exercise 115.7(5).) Show further that if

$$0 \to M' \xrightarrow{f} M \xrightarrow{g} M'' \to 0$$

is a short exact sequence of $R$-modules, then so is

$$0 \to S^{-1}R \otimes_R M' \xrightarrow{1 \otimes f} S^{-1}R \otimes_R M \xrightarrow{1 \otimes g} S^{-1}R \otimes_R M'' \to 0.$$
8. Let \( R \) be a commutative ring and \( N \) an \( R \)-module. We say that \( N \) is a flat \( R \)-module if whenever
\[
0 \to M' \xrightarrow{f} M \xrightarrow{g} M'' \to 0
\]
is exact, so is
\[
0 \to M' \otimes_R N \xrightarrow{f \otimes 1} M \otimes_R N \xrightarrow{g \otimes 1} M'' \otimes_R N \to 0.
\]
Show that if \( N \) is a projective \( R \)-module, then it is \( R \)-flat.

9. Let \( R \) be a commutative ring and \( V \) and \( W \) be finitely generated free \( R \)-modules with ordered \( R \)-bases \( B = \{v_1, \ldots, v_n\} \) and \( C = \{w_1, \ldots, w_m\} \) respectively. Let \( D = \{v_1 \otimes w_1, \ldots, v_1 \otimes w_m, v_2 \otimes w_1, \ldots, v_n \otimes w_m\} \), an ordered \( R \)-basis for \( V \otimes_R W \). Let \( T \in \text{End}_R(V) \) and \( S \in \text{End}_R(W) \). Write \( A = [T]_B \) and \( B = [S]_C \). Show that the matrix representation \( [T \circ S]_D \) is the Kronecker product
\[
\begin{pmatrix}
A_{11}B & \cdots & A_{1m}B \\
\vdots & \ddots & \vdots \\
A_{n1}B & \cdots & A_{nm}B
\end{pmatrix}
\]
i.e.,
\[
([T \circ S]_D)_{n(i-1)+r,m(j-1)+s} = ([T]_B)_{ij} ([S]_C)_{rs}.
\]


12. In Properties 116.4, which properties can be generalized to the case when the rings are not necessarily commutative? State and prove your answer.

117. Tensor, Symmetric, and Exterior Algebras

In the last section, we defined the tensor product of modules over a commutative ring. In this section, we show how to construct algebras from given algebras and given modules.

First let us recall a ring \( A \) over a commutative ring \( R \) is called an \emph{\( R \)-algebra} if there exists a ring homomorphism \( \varphi_A : R \to Z(A) \) (where \( Z(A) \) is the center of \( A \)) and an \emph{\( R \)-algebra homomorphism} of \( R \)-algebras is a ring homomorphism \( \psi : A \to B \) of \( R \)-algebras such that
\[
\begin{tikzcd}
A \arrow[swap]{dr}{\varphi_A} & B \arrow{dl}{\psi} \\
& R
\end{tikzcd}
\]
commutes. The algebra structure on \( A \) makes \( A \) into a left and right \( R \)-module. For example, every ring is a \( \mathbb{Z} \)-algebra. If \( R \) is a field \( F \) and \( A \) a nonzero \( F \)-algebra, we often identify \( F \) and \( F1_A \), i.e., view \( \varphi_A \) as an inclusion as before. An \( R \)-algebra \( A \) is called a graded \( R \)-algebra if \( A = \bigoplus_{i=0}^{\infty} A_i \) as an additive group with \( \varphi_A(R) \subset A_0 \) and \( A_iA_j \subset A_{i+j} \) for all \( i, j \). In particular, \( A_0 \) is a ring. Elements in \( A_i \) are called homogeneous elements of degree \( i \). If \( B = \bigoplus_{i \geq 0} B_i \) is another graded \( R \)-algebra, an \( R \)-algebra homomorphism \( \psi : A \to B \) is called graded if \( \varphi(A_i) \subset B_i \) for all \( i \geq 0 \). We want to construct some graded \( R \)-algebras.

Let \( R \) be a commutative ring and \( S \) a set. Then the free commutative \( R \)-algebra on \( S \) is the \( R \)-algebra satisfying the appropriate universal property — what is it? It exists and is unique up to a unique isomorphism, i.e., if two \( R \)-algebras satisfy the universal property there is a unique \( R \)-algebra isomorphism between them. Indeed, \( R[t_s] \), the polynomial ring on \( |S| \)-variables, works.

**Construction 117.1.** Let \( R \) be a commutative ring and \( A, B \) two \( R \)-algebras. If \( (a, b) \in A \times B \), then

\[
m_{a,b} : A \times B \to A \otimes_R B \text{ given by } (a', b') \mapsto aa' \otimes bb'
\]

is \( R \)-bilinear, so induces a unique \( R \)-homomorphism

\[
\mu_{a,b} : A \otimes_R B \to A \otimes_R B \text{ given by } a' \otimes b' \mapsto aa' \otimes bb'
\]

and hence

\[
(A \otimes_R B) \times (A \otimes_R B) \to A \otimes_R B \text{ given by } (a \otimes b, a' \otimes b') \mapsto aa' \otimes bb'
\]

is a well defined map that makes \( A \otimes_R B \) into a ring and into an \( R \)-algebra by \( R \to A \otimes_R B \) by \( r \mapsto r \otimes 1 \), where \( r \otimes 1 = r1_A \otimes 1_B = 1_A \otimes r1_B = 1 \otimes r = \varphi_A(r)1_A \otimes 1_B = 1_A \otimes \varphi_B(r)1_B \) for all \( r \) in \( R \). We also call this algebra the tensor product of \( A \) and \( B \).

**Definition 117.2.** Let \( R \) be a commutative ring and \( M, N \) be \( R \)-modules. For each positive integer \( n \), let \( T^n(M) \) be the \( R \)-module \( M \otimes_R \cdots \otimes_R M \). Then \( T^n(M) \) is generated by \( m_1 \otimes \cdots \otimes m_n \) with \( m_i \in M \). Define \( T^0(M) := R \). If \( f : M \to N \) is an \( R \)-homomorphism, then there exists a unique \( R \)-homomorphism \( T^n f : T^n(M) \to T^n(N) \) induced by the \((R-)n\)-linear map \( M_1 \times \cdots \times M_n \to T^n(N) \) given by \((m_1, \cdots, m_n) \mapsto f(m_1) \otimes \cdots \otimes f(m_n) \) (why?).

**Example 117.3.** Let \( n \) be a positive integer, \( V \) an \( m \)-dimensional \( F \)-vector space on basis \( B = \{v_1, \ldots, v_m\} \). Then \( T^n(V) \) is an \( F \)-vector space basis \( \{v_{i_1} \otimes \cdots \otimes v_{i_n} \mid 1 \leq i_j \leq m \} \) (repetitions are allowed), a
vector space of dimension \( m^n \). [The same is true if we have \( M \) a free \( R \)-module of on a basis \( B \) with \( R \) a commutative ring.] We also define \( T^{-n}(V) := T^n(V^*) \), with \( V^* \) the dual space \( \text{Hom}_F(V,F) \).

**Construction 117.4.** Let \( R \) be a commutative ring and \( M \) an \( R \)-module. Associativity of the tensor products implies that

\[
T^m(M) \times T^n(M) \rightarrow T^{m+n}(M) \quad \text{given by} \quad (x, y) \mapsto x \otimes y
\]

is \( R \)-bilinear, so induces a (non-commutative) ring structure on the \( R \)-module

\[
T(M) := \prod_{i=0}^{\infty} T^n(M)
\]

called the *tensor algebra* of \( M \). If \( x, y \in T(M) \), we write \( xy \) for the product. Hence if \( x \in T^m(M) \) and \( y \in T^n(M) \), then \( xy = x \otimes y \). Under this ring structure, we have

\[
0_{T(M)} := \prod_{n=0}^{\infty} 0_{T^n(M)}
\]

\[
1_{T(M)} := 1_R
\]

\[
R := T^0(M), \text{ a subring of } T(M)
\]

\( T(M) \) is a \( R \)-module

\( T(M) \) is an \( R \)-algebra.

The \( R \)-algebra \( T(M) \) also satisfies the following universal property: If \( \varphi : M \rightarrow A \) is a \( R \)-homomorphism of \( R \)-algebras, then there exists a unique \( R \)-algebra homomorphism \( \Phi : T(M) \rightarrow A \) satisfying \( \Phi|_M = \varphi \).

Let \( R \) be a commutative ring and \( M \) an \( R \)-module. As the tensor algebra \( T(M) \) is a ring, it has ideals, hence quotient rings. We use this to construct other algebras. But \( T(M) \) has an additional structure, it is a graded \( R \)-algebra by \( T(M) = \bigoplus T^n(M) \) as the ring multiplication on \( T(M) \) satisfies \( T^m(M) \times T^n(M) \rightarrow T^{m+n}(M) \). We want to construct graded \( R \)-algebras from \( T(M) \). The first is the *symmetric algebra*.

**Construction 117.5.** Let \( R \) be a commutative ring and \( M \) an \( R \)-module. For each positive integer \( n \), let \( \mathfrak{A}_n \) be the submodule of \( T^n(M) \) generated by all elements \( m_1 \otimes \cdots \otimes m_n - m_{\sigma(1)} \otimes \cdots \otimes m_{\sigma(n)} \) with
Define \( A_0 = 0 \). Let

\[
\mathfrak{A} := \bigoplus_{i \geq 0} \mathfrak{A}_n
\]

\[
S^n(M) := T^n(M)/\mathfrak{A}_n
\]

\[
S(M) := \bigoplus_{i \geq 0} S^n(M) = \bigoplus_{i \geq 0} T^n(M)/\mathfrak{A}_n.
\]

Check that \( \mathfrak{A} \) is an ideal in \( T(M) \) and \( S(M) = T(M)/\mathfrak{A} \) is a graded \( R \)-algebra, i.e., the ring multiplication on the algebra \( S(M) \) satisfies

\[
S^m(M) \times S^n(M) \to S^{m+n}(M)
\]

for all \( m, n \geq 0 \). The \( R \)-algebra \( S(M) \) is a commutative algebra called the symmetric algebra of \( M \) and we have ring homomorphisms \( R \xrightarrow{\varphi_T(M)} T(M) \xrightarrow{-} S(M) \) with \( - \) the canonical epimorphism.

**Example 117.6.** Let \( V \) be an \( n \)-dimensional \( F \)-vector space. Then \( S(V) \cong F[t_1, \ldots, t_n] \) (An analogous statement holds if \( V \) is an \( R \)-free module of rank \( n \) with \( R \) commutative):

Let \( \mathcal{B} = \{v_1, \ldots, v_n\} \) be an \( F \)-basis for \( V \) and \( \varphi : T^m(V) \to S^m(V) \) the canonical \( R \)-module epimorphism. Then

\[
\{v_{i_1} \otimes \cdots \otimes v_{i_m} \mid 1 \leq i_1 \leq \cdots \leq i_m \leq n\}
\]

is an \( F \)-basis for \( S^m(V) \). In particular,

\[
\dim_F S^m(V) = \binom{m+n-1}{n-1} = \binom{m+n-1}{m}
\]

the number of (homogeneous) monic monomials of total degree \( m \) in \( t_1, \ldots, t_n \). (Why?) In particular, \( \dim_F S(V) = \infty \).

In the construction above, the ideal \( \mathfrak{A} := \bigoplus_{n \geq 0} \mathfrak{A}_n \) satisfies

\[
\mathfrak{A} \cap T^n(M) = \mathfrak{A}_n \quad \text{and} \quad \mathfrak{A}_m \mathfrak{A}_n \subset \mathfrak{A}_{mn}
\]

for all \( m, n \geq 0 \).

We say \( \mathfrak{A} = \bigoplus_{n \geq 0} \mathfrak{A}_n \) is a graded or homogeneous ideal in the graded ring \( T(M) \).

We want universal properties for \( S^n(M) \) and \( S(M) \). Call a map \( f : X \times \cdots \times X \to Y \) of \( R \)-modules symmetric if \( f(x_1, \ldots, x_n) = f(x_{\sigma(1)}, \ldots, x_{\sigma(n)}) \) for all \( \sigma \in S_n \).

Let \( R \) be a commutative ring and \( M \) an \( R \)-module. Then we have a canonical \( n \)-linear map \( f : \underbrace{M \times \cdots \times M}_n \to S^n(M) \) given by the composition \( \underbrace{M \times \cdots \times M}_n \to T^n(M) \to T^n(M)/\mathfrak{A}_n \). Then \( S^r(M) \) satisfies the following universal property:
Let \( g : M \times \cdots \times M \to N \) be an \( n \)-linear and symmetric map of \( R \)-modules. Then there exists a unique \( R \)-homomorphism \( \varphi : S^n(M) \to N \) such that

\[
\begin{array}{c}
M \times \cdots \times M \xrightarrow{f} S^n(M) \\
\downarrow{g} \quad \downarrow{\varphi} \\
N
\end{array}
\]

commutes and the commutative \( R \)-algebra \( S(M) \) satisfies the following universal property:

If \( A \) is a commutative \( R \)-algebra and \( \varphi : M \to A \) is an \( R \)-homomorphism, then there exists a unique \( R \)-algebra homomorphism

\[
(117.8) \quad \Phi : S(M) \to A \text{ satisfying } \Phi|_M = \varphi.
\]

We next construct another quotient of the tensor algebra, this time a non-commutative algebra.

**Construction 117.9.** Let \( R \) be a commutative ring and \( M \) an \( R \)-module. Define \( \mathfrak{B}_n \) to be the submodule of \( T^n(M) \) generated by all \( m_1 \otimes \cdots \otimes m_n \) with \( m_1, \ldots, m_n \) in \( M \) and \( m_i = m_j \) for some \( i \neq j \). Set \( \wedge^n(M) := T^n(M)/\mathfrak{B}_n \). The map \( f : M \times \cdots \times M \to \wedge^n(M) \) defined by the composition \( M \times \cdots \times M \to T^n(M) \to T^n(M)/\mathfrak{B}_n \) is \( n \)-linear and alternating, i.e., satisfies \( f(x_1, \ldots, x_n) = 0 \) if there exists an \( i \neq j \) such that \( x_i = x_j \). It satisfies the following universal property:

If \( N \) is an \( R \)-module and \( g : M \times \cdots \times M \to N \) is \( n \)-linear and alternating, then there exists a unique \( R \)-homomorphism \( \varphi : \wedge^n(M) \to N \) such that

\[
\begin{array}{c}
M \times \cdots \times M \xrightarrow{f} \wedge^n(M) \\
\downarrow{g} \quad \downarrow{\varphi} \\
N
\end{array}
\]

commutes. This follows as there exists a unique \( R \)-homomorphism \( \psi : T^n(M) \to N \) such that

\[
\begin{array}{c}
M \times \cdots \times M \xrightarrow{f'} T^n(M) \\
\downarrow{g} \quad \downarrow{\psi} \\
N
\end{array}
\]
commutes. Then

\[
\begin{array}{c}
0 \rightarrow B_n \xrightarrow{h} T^n(M) \xrightarrow{-\psi} \Lambda^n(M) \rightarrow 0 \\
\downarrow 0 \quad \downarrow \psi \\
N
\end{array}
\]

is exact and commutes, so \(\varphi\) exists and is unique by the universal property of the cokernel.

We denote the image of \((m_1, \ldots, m_n)\) in \(\Lambda^n(M)\) by \(m_1 \wedge \cdots \wedge m_n\).

If \(\varphi : M \rightarrow N\) is an \(R\)-homomorphism, then

\[
\underbrace{M \times \cdots \times M}_n \rightarrow \Lambda^n(N)
\]

defined by \(m_1, \ldots, m_n \mapsto \varphi(m_1) \wedge \cdots \wedge \varphi(m_n)\) is \(n\)-linear and alternating, so induces a unique \(R\)-homomorphism \(\wedge^n \varphi : \Lambda^n(M) \rightarrow \Lambda^n(N)\).

Set \(\Lambda^0(M) := R, \Lambda^1(M) = M\) and \(\Lambda(M) := \bigoplus_{n \geq 0} \Lambda^n(M)\).

**Check.** \(\mathcal{B} := \bigoplus_{n \geq 0} B_n\) is a graded ideal in \(T(M)\) by using

\[
T^m(M) \times \mathcal{B}_n \quad \text{and} \quad \mathcal{B}_m \times T^n(M) \quad \text{lie in} \quad \mathcal{B}_{m+n},
\]

so induces an \(R\)-bilinear map \(\Lambda^m(M) \times \Lambda^n(M) \rightarrow \Lambda^{m+n}(M)\) and (as we have associativity) an \(R\)-homomorphism

\[
\Lambda^m(M) \otimes_R \Lambda^n(M) \rightarrow \Lambda^{m+n}(M).
\]

This defines a graded \(R\)-algebra structure \(\Lambda(M)\) with multiplication \(\Lambda(M) \times \Lambda(M) \rightarrow \Lambda(M)\) defined by

\[(x_1 \wedge \cdots \wedge x_m) \cdot (y_1 \wedge \cdots \wedge y_n) := x_1 \wedge \cdots \wedge x_m \wedge y_1 \wedge \cdots \wedge y_n.\]

[Why does it suffice to define multiplication on generators?] This \(R\)-algebra is called the **exterior algebra** of \(M\).

**Remark 117.10.** Let \(R\) be a commutative ring and \(M\) an \(R\)-module.

1. \(\Lambda(M)\) can be defined to be the \(R\)-algebra generated by \(M\) with defining relations

\[m \wedge m = 0 \quad \text{for all} \quad m \in M:\]

If \(m, n\) are elements of \(M\), then

\[mn + nm = (n + m)^2 - n^2 - m^2 \quad \text{in} \quad T(M),\]

so

\[m \wedge n = -n \wedge m \quad \text{in} \quad \Lambda(M)\]

and it follows that the graded ideal \(\mathcal{B}\) is generated in the algebra \(T(M)\) by \(m^2, m \in M\).
2. Using (1), one checks \( \wedge(M) \) satisfies the following universal property: If \( A \) is an \( R \)-algebra and \( \varphi : M \to A \) an \( R \)-homomorphism satisfying \( \varphi(m)^2 = 0 \) for all \( m \in M \), then there exists a unique \( R \)-algebra homomorphism \( \Phi : \wedge(M) \to A \) satisfying \( \Phi|_M = \varphi \).

In particular, any \( R \)-homomorphism \( \varphi : M \to N \) extends to a unique graded \( R \)-algebra homomorphism \( \wedge \varphi : \wedge(M) \to \wedge(N) \), i.e., \( \wedge \varphi(\wedge^n(M)) \subset \wedge^n(N) \). For example, \( \varphi(1_M) = 1_{\wedge(M)} \).

3. If \( \varphi : M \to N \) and \( \psi : N \to P \) are \( R \)-homomorphisms, then \( \wedge(\psi \varphi) = \wedge \psi \circ \wedge \varphi \).

**Exercises 117.11.**

1. Let \( V \) be an \( n \)-dimensional \( F \)-vector space. Show

\[
\dim_F S^m(V) = \binom{m + n - 1}{n - 1} = \binom{m + n - 1}{m}.
\]

2. Let \( R \) be a commutative ring. Define the free commutative \( R \)-algebra on a set \( S \) and show it is a polynomial ring on \( |S| \) variables.

3. Let \( R \) be a commutative ring. Define the free \( R \)-algebra on a set \( S \) and show it exists.

4. Let \( R \) be a commutative ring and \( A \) and \( B \) commutative \( R \)-algebras. Show that \( A \otimes_R B \) is the coproduct of \( A \) and \( B \) as rings [definition?]. [If \( A \) and \( B \) are not commutative, this is not true.]

5. Let \( K_i/F \) be (arbitrary) field extensions for \( i = 1, 2 \). A composite of \( K_1, K_2 \) over \( F \) is a field \( L \) such that there exist (field) homomorphisms \( \varphi_i : K_i \to L \) fixing \( F \), i.e., \( F \)-homomorphisms, such that \( L \) is generated by \( \varphi_1(K_1) \varphi_2(K_2) \). Show that there is a bijection between \( \text{Spec}(K_1 \otimes_F K_2) \) and \( F \)-isomorphism classes of composites of \( K_1 \) and \( K_2 \) over \( F \).

6. If \( R \) is a commutative ring and \( M \) an \( R \)-module, prove that the tensor algebra \( T(M) \) satisfies the universal property in Construction 117.4.

7. Fill in the details needed to justify the construction of the symmetric algebra in Construction 117.5.

8. Prove \( S^r(M) \) satisfies the universal property (117.7).

9. Prove \( S(M) \) satisfies the universal property (117.8).

10. Prove \( \wedge(M) \) satisfies the universal property in Remark 117.10(2).

**118. The Determinant**

In this section, we show that the determinant exists for matrices over a commutative ring using the exterior algebra of a finitely
generated free module $M$. We construct the determinant of an $R$-endomorphism of $M$ that gives rise to the determinant of its matrix representations, so is more intrinsic than the more elementary matrix proof. Recall over a commutative ring all bases for a finitely generated free module have have the same number of elements (the rank of the module). Unlike the tensor algebra over a finitely free module that has infinite rank, the exterior algebra has finite rank. More precisely, we have the following:

**Theorem 118.1.** Let $R$ be a commutative ring and $M$ a free $R$-module of rank $n$ on basis $\{x_1, \ldots, x_n\}$. Then $\bigwedge(M)$ is a free $R$-module of rank $2^n$. More precisely, 

$$B_r := \{x_{i_1} \wedge \cdots \wedge x_{i_r} \mid 1 \leq i_1 < \cdots < i_r \leq n\}$$

is a basis for $\bigwedge^r(M)$ for $r \leq n$. In particular,

$$\operatorname{rank}\bigwedge^r(M) = \begin{cases} \binom{n}{r} & \text{if } r \leq n \\ 0 & \text{otherwise.} \end{cases}$$

To prove this we need a preliminary idea and its development in the case of exterior algebras.

**Construction 118.2.** Let $R$ be a commutative ring, $M$ an $R$-module, and $A = \prod_{i \geq 0} A_i$ a graded $R$-algebra.

For each $i \geq 0$, define an $R$-homomorphism $\sigma : A_i \to A_i$ by $z \mapsto (-1)^iz$. Then $\sigma$ induces a graded $R$-algebra homomorphism $\sigma : A \to A$, since $(-1)^i(-1)^j = (-1)^{i+j}$. Note that $\sigma^2 = 1_A$.

An $R$-homomorphism $\delta : A \to A$ is called an anti-derivation if $\delta(ab) = \delta(a)b + \sigma(a)\delta(b)$ for all $a, b \in A$.

If $A$ is $T(M)$, $S(M)$, or $\bigwedge(M)$, then $M$ generates $A$ as an $R$-algebra, hence $\delta$ is determined by what it does to $A_1 = M$.

**Case.** $A = T(M)$:

Let $\tilde{\varphi} : M \to T(M)$ be an $R$-homomorphism. Then $M \times M \to T(M)$ defined by

$$(*) \quad (a, b) \mapsto \tilde{\varphi}(a)b - a\tilde{\varphi}(b) = \tilde{\varphi}(a)b + \sigma(a)\tilde{\varphi}(b)$$

is $R$-bilinear, so induces an $R$-homomorphism $T^2(M) \to T(M)$. Inductively, we obtain $R$-homomorphisms $T^n(M) \to T(M)$, i.e., $\tilde{\varphi}$ induces an anti-derivation $T(M) \to T(M)$. We also call this map $\tilde{\varphi}$.

**Case.** $A = \bigwedge(M)$:
Let $\varphi$ be the composition

$$T(M) \xrightarrow{\tilde{\varphi}} T(M) \rightarrow \wedge(M),$$

with $\tilde{\varphi}$ the anti-derivation of the previous case.

**Claim.** If $\varphi$ satisfies

$$(118.3) \quad \varphi(x) \wedge x - x \wedge \varphi(x) = 0 \text{ for all } x \in M,$$

then $\varphi$ induces an anti-derivation $\wedge(M) \rightarrow \wedge(M)$ that we shall also call $\varphi$:

As $\tilde{\varphi} : T(M) \rightarrow T(M)$ is an anti-derivation, to show that $\varphi$ induces an anti-derivation $\wedge(M) \rightarrow \wedge(M)$, it suffices to show that

$$\tilde{\varphi}(\mathfrak{B}) \subset \mathfrak{B} = \ker(- : T(M) \rightarrow \wedge(M)).$$

Let $x \in M$. By (*), we have

$$(\dagger) \quad \tilde{\varphi}(x^2) = \tilde{\varphi}(x \cdot x) = \tilde{\varphi}(x)x + \sigma(x)\tilde{\varphi}(x) = \tilde{\varphi}(x)x - x\tilde{\varphi}(x).$$

Applying the map $-$ to Equation $(\dagger)$, shows that $\tilde{\varphi}(x^2)$ lies in $\mathfrak{B}$. Let $a$ and $b$ lie in $T(M)$ and $x$ in $M$. Then

$$\tilde{\varphi}(ax^2b) = \tilde{\varphi}(a)x^2b + \sigma(a)\tilde{\varphi}(x^2b)$$

$$= \tilde{\varphi}(a)x^2b + \sigma(a)\tilde{\varphi}(x^2)b + \sigma(a)\sigma(x^2)\tilde{\varphi}(b)$$

$$= \tilde{\varphi}(a)x^2b + \sigma(a)\tilde{\varphi}(x^2)b + \sigma(a)\sigma(x)^2\tilde{\varphi}(b).$$

Since $x^2, \sigma(x)^2$ lie in $\mathfrak{B}$ and $\tilde{\varphi}(x^2)$ lies in $\mathfrak{B}$ by $(\dagger)$ with $\mathfrak{B}$ a 2-sided ideal in $T(M)$, we conclude that $\tilde{\varphi}(ax^2b)$ lies in $\mathfrak{B}$ for all $a, b$ in $T(M)$ and for all $x$ in $M$. Since $\mathfrak{B}$ is generated by $x^2, x \in M$, by Remark 117.10(1), we have

$$\mathfrak{B} = \langle ax^2b \mid a, b \in T(M), x \in M \rangle.$$ 

Therefore, $\tilde{\varphi}(\mathfrak{B}) \subset \mathfrak{B}$, and the claim is established.

**Examples 118.4.** Let $R$ be a commutative ring and $M$ an $R$-module. Suppose that we are in the case $A = \wedge(M)$ above together with its notation.

1. If $\varphi(x)$ lies in $M$ for all $x$ in $M$, then equation (118.3) becomes

$$2\varphi(x) \wedge x = 0 \text{ for all } x \in M.$$

2. If $\varphi(x)$ lies in $R$ for all $x$ in $M$, then equation (118.3) always holds. In particular, as $R = \wedge^0(M) \subset \wedge(M)$, if $\varphi$ lies in $M^* = \text{Hom}_R(M, R)$, then $\varphi$ induces an anti-derivation $\wedge(M) \rightarrow \wedge(M)$ that we also call $\varphi$. 

3. Suppose that $M$ is $R$-free on basis $\mathcal{B} = \{x_1, \ldots, x_n\}$ and $B^* = \{f_1, \ldots, f_n\}$, the dual basis for $M^* = \text{Hom}_R(M, R)$. Then each $f_i$ defines an anti-derivation $f_i : \Lambda(M) \rightarrow \Lambda(M)$.

We now proceed to the proof of Theorem 118.1.

**Proof.** We use the notation of Example 118.4(3). As $x_{\sigma(i_1)} \wedge \cdots \wedge x_{\sigma(i_r)} = \pm x_{i_1} \wedge \cdots \wedge x_{i_r}$, for all $\sigma \in S_r$, we have $B_r = \{x_{i_1} \wedge \cdots \wedge x_{i_r} \mid 1 \leq i_1 < \cdots < i_r \leq n\}$ generates $\Lambda^r(M)$. As $\Lambda(M) = \bigoplus_{r \geq 0} \Lambda^r(M)$ is graded, it suffices to show that $B_r$ is $R$-linearly independent. Suppose this is false. Since $\mathcal{B} = \mathcal{B}_1$ is linearly independent, there exists a minimal $r$ such that $B_r$ is linearly dependent. Let

(*)\[ \sum_{1 \leq i_1 < \cdots < i_r \leq n} a_{i_1, \ldots, i_r} x_{i_1} \wedge \cdots \wedge x_{i_r} = 0 \] in $\Lambda^r(M)$, not all $a_{i_1, \ldots, i_r}$ zero. Changing notation, we may assume that $a_{j_1, j_2, \ldots, j_r}$ is nonzero with $1 < j_2 < \cdots < j_r \leq n$. Applying the anti-derivation arising from $f_1$ in $\mathcal{B}^*$ to (*), we get

\[ \sum_{1 \leq i_2 < \cdots < i_r \leq n} a_{i_1, i_2, \ldots, i_r} x_{i_2} \wedge \cdots \wedge x_{i_r} = 0, \]

since $f_1(x_i) = \delta_{i1}$. This implies that $B_{r-1}$ is linearly dependent, a contradiction. \hfill \Box

**Corollary 118.5.** The sign map $\text{sgn} : S_n \rightarrow \{\pm 1\}$ given by

\[ \text{sgn} \, \sigma = \begin{cases} 1 & \text{if } \sigma \text{ is a product of an even number of transpositions} \\ -1 & \text{otherwise} \end{cases} \]

is a well-defined group homomorphism.

**Proof.** Let $R$ be a commutative ring and $M$ a free $R$-module on basis $\{x_1, \ldots, x_n\}$, then $\Lambda^n(M)$ is $R$-free on basis $\{x_1 \wedge \cdots \wedge x_n\}$. Hence if $\sigma \in S_n$, there exists a unique $\text{sgn}(\sigma)$ in $\{\pm 1\}$ satisfying

(118.6)\[ x_{\sigma(1)} \wedge \cdots \wedge x_{\sigma(n)} = \text{sgn}(\sigma)x_1 \wedge \cdots \wedge x_n. \]

It is easily checked that $\text{sgn}$ is a group homomorphism. \hfill \Box

**Construction 118.7.** Let $R$ be a commutative ring and $M$ an $R$-free module on basis $\{x_1, \ldots, x_n\}$. Let $f : M \rightarrow M$ be an $R$-homomorphism. Write

\[ f(x_i) = \sum_{j=1}^{n} \alpha_{ji} x_j \quad \text{for } i = 1, \ldots, n \]
and $A = (\alpha_{ij}) = [f]_B$, the matrix representation of $f$ relative to the (ordered) basis $B$. Then $f$ induces a unique graded $R$-algebra homomorphism $\wedge f : \Lambda(M) \to \Lambda(M)$. Since $\{x_1 \wedge \cdots \wedge x_n\}$ is a basis for $\Lambda^n(M)$, there exists a unique element $\lambda$ in $R$ satisfying $\wedge^n f = \lambda 1_{\Lambda^n(M)}$, where $\wedge f|_{\Lambda^n(M)} = \wedge^n f$. In particular, $\lambda$ is an eigenvalue of $\wedge^n f$ on the rank one free $R$-module $\Lambda^n(M)$. Using equation (118.6) and the alternating property, we see that

$$\lambda 1_{\Lambda^n(M)} = \wedge^n f(x_1 \wedge \cdots \wedge x_n)$$

$$= f(x_1) \wedge \cdots \wedge f(x_n) = \sum_{j=1}^n \alpha_{j1} x_j \wedge \cdots \wedge \sum_{j=1}^n \alpha_{jn} x_j$$

$$= \sum_{S_n} \text{sgn}(\sigma) \alpha_{\sigma(1)1} \cdots \alpha_{\sigma(n)n} x_1 \wedge \cdots \wedge x_n.$$ 

The element

$$\lambda = \sum_{S_n} \text{sgn}(\sigma) \alpha_{\sigma(1)1} \cdots \alpha_{\sigma(n)n}$$

is called the determinant of $f$ and denoted $\det f$. As $\det f$ is an eigenvalue for $\wedge^n f$ on a rank one free $R$-module, it is unique and independent of $B$. The map $\det : \text{End}_R(M) \to R$ defined by $\det f = \wedge^n f$ is $n$-multilinear, alternating and satisfies $\det 1_V = 1_R$. It is called the determinant (function).

**Remark 118.8.** The above shows the existence and uniqueness of the usual determinant (function) on matrices $\mathbb{M}_n(R)$ over a commutative ring $R$, i.e., it is the unique alternating, $n$-multilinear function $D : \mathbb{M}_n(R) \to R$ on columns (or rows) of $n \times n$ matrices over $R$ satisfying $D(I) = 1$. Indeed, if $\mathcal{S}$ is the standard basis for $R^n$, then $D(A) = \det(A)$, viewing $A$ as a linear transformation on $R^n$. In particular, $\det(A) = 0$ if and only if $A$ is not a unit in $\mathbb{M}_n(R)$.

Next let $g : M \to M$ be another $R$-homomorphism and $B = [g]_B$. Then

$$\wedge^n(fg) = \wedge^n f \circ \wedge^n g \text{ so } \det fg = \det f \det g, \text{ i.e., } \det AB = \det A \det B.$$ 

More generally, $f, g$ induce $\wedge^r f, \wedge^r g : \Lambda^r(M) \to \Lambda^r(M)$ for each $r = 1, \ldots, n$. Set $A^{(r)} = [\wedge^r f]_{B_r}$ and $B^{(r)} = [\wedge^r g]_{B_r}$ in $\mathbb{M}_n(R)$. Then

$$(AB)^{(r)} = A^{(r)} B^{(r)}$$

which gives identities of degree $r$ in the entries of $A$ and $B$ called the Binet-Cauchy Equations. Note that $A^{(r)}$ is an $\left( \binom{n}{r} \times \binom{n}{r} \right)$-matrix whose entries are the $r$th order minors of $A$, etc.
To get additional formulae, let \( A = (a_{ij}) \) be an \( n \times n \)-matrix over a commutative ring \( R \). Then we view \( A \) as the matrix representation of the appropriate \( f : R^n \rightarrow R^n \) in the standard basis \( \{ e_1, \ldots, e_n \} \). Set \( S = \{ 1, \ldots, n \} \). If \( H \) is a subset of \( S \) consisting of \( r \) elements, let \( H \) be the ordered \( r \)-tuple \( (i_1, \ldots, i_r) \) where \( 1 \leq i_1 < \cdots < i_r \leq n \) with all \( i_j \in H \). and \( e_H = e_{i_1} \land \cdots \land e_{i_r} \). If \( K \) is another subset of \( S \), then

\[
(118.9) \quad e_H \land e_K = \begin{cases} \varepsilon_{H,K} e_{H \cup K} & \text{if } H \cap K = \emptyset \\ 0 & \text{if } H \cap K \neq \emptyset, \end{cases}
\]

where \( \varepsilon_{H,K} = (-1)^m \), if \( H \cap K = \emptyset \) and \( m \) is the number of transpositions modulo two to get \( H \cup K \) from \( H \cup K \). If \( H \) is a subset of \( S \), we denote by \( H' \) the complement of \( H \) in \( S \). Fix \( r, 1 \leq r \leq n \), and let \( S(r) = \{ I \mid I \in P(S) \text{ with } |I| = r \} \), where \( P(S) \) is the power set of \( S \), i.e., the set of subsets of \( S \). Let \( A_{L,H} \) denote the submatrix of \( A \) with rows of \( A \) associated to \( L, I \in S(r) \), and columns of \( A \) associated to \( H \) i.e., \( \det A_{L,H} \) is the \( r \)-order minor associated to \( L' \) and \( H' \).

Fix \( H \in S(r) \). Then check we have

\[
(118.10) \quad \land^r (f)(e_H) = \sum_{I \in S(r)} \det(A_{L,H}) e_L.
\]

Therefore, for such a fixed \( H \), we have

\[
(\dagger) \quad (\det A \cdot \varepsilon_{H,H'}) e_S = (\land^r (f))(e_H) \land (\land^{n-r} (f))(e_{H'})
\]

\[
= \sum_{I \in S(r)} \sum_{J \in S(n-r)} \det(A_{L,H}) \det(A_{L',H'}) e_L \land e_J
\]

\[
= \sum_{I \in S(r)} \sum_{J \in S(n-r)} \det(A_{L,H}) \det(A_{L',H'}) \varepsilon_{I,J} e_S
\]

\[
= \sum_{I \in S(r)} \varepsilon_{I,L'} \det(A_{L,H}) \det(A_{L',H'}) e_S
\]

by (118.9). In particular,

\[
\det A = \sum_{I \in S(r)} \varepsilon_{H,H'} \varepsilon_{I,L'} \det(A_{L,H}) \det(A_{L',H'}).
\]

This is called the Laplace expansion of \( \det A \). For \( r = 1 \), this is the well-known expansion by minors along rows. By symmetry, we have an analogous result on columns by using right modules. If we let \( A(i,j) \) be the \((n-1)\)-order minor of \( A \), i.e., the determinant of the submatrix of \( A \) by deleting the \( i \)th row and \( j \)th column, then \((-1)^{i+j} A(i,j)\) is called the \((i,j)\)th cofactor of \( A \). Define the classical adjoint \( \text{adj}(A) \) of
A to be the transpose of the matrix of cofactors of A, i.e., the \((i, j)\) entry of \(\text{adj}(A)\) is \((-1)^{i+j}A(j, i)\). By Laplace expansion, we check that

\[
\sum_{i=1}^{n} a_{ij} (-1)^{i+1} A(j, 1) + \cdots + a_{in} (-1)^{i+1} A(j, n)
\]

is the expansion of a determinant whose \(j\)th row equals the \(i\)th row of \(A\) while the other rows are as in \(A\). If \(i \neq j\) this means this matrix has two rows the same so is zero, hence is zero off the main diagonal. It follows that

\[
A \cdot \text{adj}(A) = \det(A) \cdot I
\]

with \(I\) the \(n \times n\) identity matrix. Similarly, \(\text{adj}(A) \cdot A = \det(A) \cdot I\).

Let \(M = \sum_{i=1}^{n} R x_i\) be an \(R\)-module. If we have a system of equations \(\sum_{j=1}^{n} a_{ji} x_j = 0\) for \(i = 1, \ldots, n\), then we must have \(\det(A)M = 0\). In particular, \(\det A\) is in the annihilator of \(M\). It follows that if \(\text{ann}_R(M) = 0\) and \(M \neq 0\), then \(\det(A) = 0\). Next suppose that \(A\) is invertible, i.e., \(\det(A)\) is a unit (respectively, \(R\) is a domain and \(M = R\)) and we want to solve a system of equations \(\sum_{j=1}^{n} a_{ji} x_j = b_i\) for \(i = 1, \ldots, n\) in \(M\) (respectively the quotient field of \(R\)), i.e., if \(x = (x_1 \cdots x_n)^t\) and \(b = (b_1 \cdots b_n)^t\) column matrices, then we want to solve the matrix equation \(Ax = b\). Then the solution is \(x = (\det A)^{-1} \text{adj}(A)b\). This is called Cramer’s Rule.

**Exercises 118.11.**

1. Let \(V\) be an \(n\)-dimensional \(F\)-vector space. If \(y_1, \ldots, y_n\) are elements in \(V\), show \(\{y_1, \ldots, y_n\}\) is linearly dependent if and only if \(y_1 \wedge \cdots \wedge y_n = 0\).

2. Let \(R\) be a commutative ring and

\[
0 \to M' \to M \to M'' \to 0
\]

an exact sequence of free \(R\)-modules of ranks \(r, n, s\), respectively. Show that there exists a natural isomorphism

\[
\varphi : \bigwedge^r (M') \otimes_R \bigwedge^s (M'') \to \bigwedge^n (M).
\]

3. Let \(M = M' \oplus M''\) be a direct sum of free \(R\)-modules of finite rank. Then for all positive integers \(n\), show that

\[
\bigwedge^n (M) \cong \prod_{r+s=n} \bigwedge^r (M') \otimes_R \bigwedge^s (M'').
\]

4. Show Equation 118.10 holds.

5. Show Equation (†) holds (in the equation right after Equation 118.10).

6. Check that Cramer’s Rule in the text is the usual one.
7. Let $R$ be a commutative ring and $A$ an $R$-algebra. A derivation $\delta : A \to A$ is an $R$-homomorphism satisfying $\delta(ab) = \delta(a)b + a\delta(b)$ for all $a, b \in A$. (E.g., the derivative on rings of real differentiable functions.) Show that if $M$ is an $R$-module, then any $R$-homomorphism $\varphi : M \to S(M)$ extends to a derivation $S(M) \to S(M)$.

8. Let $R$ be a commutative ring, $A \in \mathbb{R}^{m \times n}$, and $B \in \mathbb{R}^{n \times m}$ with $m \leq n$. Set $S = \{1, \ldots, n\}$ and $S(m) = \{I \mid I \in P(S) \text{ with } |I| = m\}$. Let $M$ be the ordered tuple $(1, 2, \ldots, m)$ where $M = \{1, \ldots, m\} \subset S$. For each $I \in S(m)$, let $A_{M,I}$ denote the $m \times m$ matrix formed from $A$ using the columns of $A$ that are elements of $I$, and $B_{I,M} \in \mathbb{M}_m(R)$ denote the $m \times m$ matrix formed from $B$ using the rows $I$ of $B$. Prove

$$\det(AB) = \sum_{I \in S(m)} \det(A_{M,I}) \det(B_{I,M}).$$

This is also called the Binet-Cauchy Equation.

9. Let $V$ be a finite dimensional real inner product space with inner product $\langle \ , \ \rangle$. Show that if $v = (v_1, \ldots, v_n)$ and $w = (w_1, \ldots, w_n)$ are elements of $V$, then

$$\sum_{1 \leq i < j \leq n} \det\begin{pmatrix} v_i & v_j \\ w_i & w_j \end{pmatrix}^2 = \det\begin{pmatrix} \langle v, v \rangle & \langle v, w \rangle \\ \langle v, w \rangle & \langle w, w \rangle \end{pmatrix}.$$

In particular, the Cauchy-Schwarz Inequality holds.

10. Let $R$ be a commutative ring. Suppose that $E \in M_n(R)$ has the form

$$E = \begin{pmatrix} A & B \\ C & D \end{pmatrix},$$

with $A \in GL_r(R)$, $B \in R^{r \times (n-r)}$, $C \in R^{(n-r) \times r}$, and $D \in \mathbb{M}_{n-r}(R)$. Let $D' = D + CA^{-1}B \in \mathbb{M}_{n-r}(R)$, called the Schur complement of $A$ in $E$ and denoted by $E/A$. Show that

$$\det E = \det A \det(E/A).$$

In particular, if $C = 0$, then $\det E = \det A \det D$.

11. Let $R$ be a commutative ring, $A \in R^{m \times n}$ and $B \in R^{n \times m}$. Show that $\det(I_m + AB) = \det(I_n + BA)$, where $I_m$ and $I_n$ are the appropriate identity matrices.

12. Let $R$ be a commutative ring, $M \in R^{n \times (n-r)}$, and $a_1, \ldots, a_{2r} \in R^{n \times 1}$. Show

$$\sum_{\sigma} (-1)^{k_1 + \cdots + k_r} \det\left( M a_{k_1} \cdots a_{k_r} \right) \det\left( M a_{k_{r+1}} \cdots a_{k_{2r}} \right) = 0,$$

where $k_1, \cdots, k_{2r}$ is a permutation of $\{1, 2, \ldots, 2r\}$ and $\sigma \in S_{2r}$ is the permutation with $1 \leq k_1 < \cdots < k_r \leq 2r$ and $k_{r+1} < \cdots < k_{2r}$.
Appendices
APPENDIX A

Axiom of Choice and Zorn’s Lemma

We have alluded to the Axiom of Choice in the text, but have used Zorn’s Lemma 27.5 for the most part. These two axioms of set theory are equivalent. In this appendix, we show this. We shall also show that Zorn’s Lemma is also equivalent to the Well-ordering Principle.

Recall from Section 27.5 the following: A pair \((S, \leq)\) (or \(S\) if \(\leq\) is clear) is called a partially ordered set or poset (under \(\leq\)) if the following holds: For all \(a, b, c \in S\), we have \(a \leq a\); if \(a \leq b\) and \(b \leq a\), then \(a = b\); and if \(a \leq b\) and \(b \leq c\), then \(a \leq c\). Let \((S, \leq)\) be a poset. We call \((S, \leq)\) a chain if for all \(a\) and \(b\) in \(S\), we have either \(a \leq b\) or \(b \leq a\). Let \(T\) be a subset of \(S\). An element \(a\) in \(S\) is called an upper bound of \(T\) if \(x \leq a\) for all \(x \in T\), and \(s\) is called a maximal element of \(S\) if \(s \leq y\) with \(y \in S\) implies that \(s = y\). We say that the poset \(S\) is inductive if every chain in \(S\) has an upper bound in \(S\). We also shall write \(a < b\) for \(a \leq b\) with \(a \neq b\). We shall also write \(\geq\) and \(>\) when it enhances the situation.

[Note: In this appendix for clarity, we write \(\subseteq\) instead of \(\subset\) as in the text.]

Zorn’s Lemma is the axiom that we assume to always be true. It says:

**Lemma A.1.** (Zorn’s Lemma) Let \(S\) be a non-empty inductive poset. Then \(S\) contains a maximal element.

Call an inductive poset \(S\) strongly inductive if every chain \(C\) in \(S\) has a least upper bound, i.e., an upper bound \(a \in S\) such that \(a \leq x\) for every upper bound \(x\) in \(S\) of \(C\). If the least upper bound for a chain lies in the chain, we shall call it a last element of the chain. We have similar definitions for lower bounds, least lower bounds, and first elements. As every strongly inductive poset is inductive, we have the following that be called a strong form of Zorn’s Lemma.

**Lemma A.2.** Assume that Zorn’s Lemma holds. Let \(S\) be a non-empty poset. If \(S\) is strongly inductive, then \(S\) contains a maximal element.

The converse is also true.
Proposition A.3. Every non-empty strongly inductive poset has a maximal element if and only if Zorn's Lemma is valid.

Proof. Suppose that $S$ is a partially ordered set and every chain in $S$ has an upper bound. We must show that $S$ has a maximal element. Let $\mathcal{P}(S) = \{A \mid A \subseteq S\}$ be the power set of $S$. It is a poset via set inclusion $\subseteq$. Let $\mathcal{C} = \{C \mid C$ is a chain in $S\}$. Since $\mathcal{C} \subseteq \mathcal{P}(S)$, it is also partially ordered by $\subseteq$. Let $C \in \mathcal{C}$ and $C_0 = \bigcup_{C \in \mathcal{C}} C$. Then $C_0$ is a least upper bound for $\mathcal{C}$. Therefore, $\mathcal{C}$ is strongly inductive, so has a maximal element $C'$ in $\mathcal{C}$. It follows that there is an upper bound $x_0$ for $C'$, viz., the union of the sets in $C'$ and this $x_0 \in C_0$, so is a maximal element of $S$. Therefore, Zorn’s Lemma holds. □

We generalize the notion of the Well-ordering Principle for the integers.

Definition A.4. Let $(S, \leq)$ be a poset. We say that $(S, \leq)$ is well-ordered if every non-empty subset $A$ of $S$ has a first element. This is equivalent to every nonempty subset of a chain in $S$ has a first element.

Definition A.5. The Well-ordering Axiom says if $S$ is a nonempty set, then there exists a partial ordering $\leq$ on $S$ such that $(S, \leq)$ is well-ordered.

Proposition A.6. Suppose every non-empty strongly inductive poset has a maximal element. Then the Well-ordering Axiom holds.

Proof. Let $S$ be a non-empty set. We must show that $S$ can be well-ordered. We know that any partial ordering on $S$ is inductive by Zorn’s Lemma, hence strongly inductive. For each $A \subseteq S$ that can be well ordered, let $W_A$ be the set of partial orderings on $A$ that are well-orderings and set $W = \{(A, \leq_A) \mid \emptyset \neq A \subseteq S, \leq_A \in W_A\}$. Clearly, $W$ is not the empty set as $(\{a\}, =)$ lies in $W$ for $a \in A$. Define a partial ordering $\leq$ on $W$ by setting $(A, \leq_A) \leq (B, \leq_B)$ in $W$, if $A \subseteq B$, $\leq_B |_{A} = \leq_A$, and if $a \in A$, $b \in B$, then $a \leq_B b$.

Let $C$ be a chain in $W$. Set $C_0 = \bigcup_{(C, \leq) \in C} C$ and define $\leq_{C_0}$ by $a \leq_{C_0} b$ if there exists $(C, \leq_C) \in C$ with $a, b \in C$ and $a \leq_C b$. As $C$ is a chain this is independent of $C$ with $a, b \in C$. Therefore, $C_0$ is well-ordered by $\leq_{C_0}$ and $(C_0, \leq_{C_0})$ is the least upper bound for $C$. Since $(C_0, \leq_{C_0})$ is strongly inductive, there exist a maximal element $(A, \leq_A)$ for $W$. If $(A, \leq_A) < W$, let $b \in W \setminus A$ and $B = A \cup \{b\}$. Define $\leq_B$ by $x \leq_B b$ for all $a \in A$ and $\leq_B |_{A} = \leq_A$. Then $(B, \leq_B) \in W$ and $(A, \leq_A) < (B, \leq_B)$. This contradicts the maximality of $(A, \leq_A)$. Hence $W = S$, and the Well-ordering Axiom holds. □
There are a number of variants for the statement of the Axiom of Choice. We show:

**Proposition A.7.** The following are equivalent:

1. If $S$ is a non-empty collection of non-empty disjoint sets $A$, then there exists a set $B$ whose elements consist of precisely one element from each set $A \in S$.

2. If $S$ is a non-empty collection of non-empty sets $A$, then there exists a function $\varphi : S \to \bigcup S A$ satisfying $\varphi(A) \in A$ for all $A \in S$.

3. If $I$ is a non-empty set and $A_i$ is a non-empty set for each $i \in I$, then the cartesian product $\prod_{i} A_i : \{ f : I \to \bigcup_{i} A_i \mid f(i) \in A_i \text{ for all } i \in I \}$ is non-empty.

**Proof.** (1) $\Rightarrow$ (3): Let $B_i = \{(i, x) \mid x \in A_i\}$, a non-empty set for each $i \in I$. Set $S = \{B_i \mid i \in I\}$, a non-empty collection of disjoint sets. By (1) there exists a set $B = \{(i, a_i) \mid \text{precisely one element } a_i \in A_i \text{ for each } i \in I\}$. Define $\varphi : I \to \bigcup_{i} A_i$ by $i \mapsto a_i$. Then $\varphi$ works.

(3) $\Rightarrow$ (2): By (3), the set $\prod_{i} S A$ is non-empty. Therefore, there exists a function $\varphi : S \to \bigcup_{i} S$ such that $\varphi(A) \in A$ for all $A \in S$.

(2) $\Rightarrow$ (1): Let $S$ be a non-empty collection of non-empty disjoint sets. By (2) there exist a function $\varphi : S \to \bigcup_{i} S S$ satisfying $\varphi(A) \in A$ for all $A \in S$. Then the image of $\varphi$ has precisely one element from each $A \in S$.

**Definition A.8.** We say the Axiom of Choice holds if the equivalent conditions of Proposition A.7 hold for all non-empty sets.

**Proposition A.9.** Suppose the Well-ordering Axiom holds. Then the Axiom of Choice holds.

**Proof.** Let $S = \{A \mid A\text{ is a non-empty set}\}$ and set $B = \bigcup S A$. As the Well-ordering Axiom holds, there exists a well-ordering $\leq_B$ on $B$. In particular, each $A \in S$ is a non-empty subset of the well-ordered set $B$. For each $A \in S$, let $a$ be the first element of $A$ under $\leq_B$. Then $\varphi : S \to \bigcup S A$ by $A \mapsto a$ satisfies $\varphi(A) \in A$ for all $A \in S$. This establishes (2) of Proposition A.7.

**Lemma A.10.** Suppose that $S$ is a strongly inductively ordered set and $\varphi : S \to S$ a function satisfying $\varphi(x) \leq x$ for every $x \in S$. If $a \in S$, then there exists an $x_0 \in S$ satisfying $x_0 \geq a$ and $\varphi(x_0) = x_0$.

**Proof.** Let $B_a$ be the collection of all subsets $B$ in $S$ that satisfy all of the following:
(i) \(a \in B\).
(ii) \(\varphi(B) \subseteq B\).
(iii) If \(C\) is a chain in \(B\), then the least upper bound of \(C\) lies in \(B\).

Since the set \(\{x \mid x \geq a\}\) lies in \(B_a\). \(B_a\) is a non-empty set. Let \(A = \bigcap_{B_a} B\). Then check that \(A \in B_a\). In particular, if \(z \in A\), then \(z \geq a\).

**Claim.** \(A\) is a chain.

Suppose that we show the claim is valid. Let \(x_0\) be the least upper bound for \(A\). Then \(x_0\) lies in \(A\), since \(S\) is strongly inductive and \(\varphi(x_0)\) lies in \(B\) by (iii). Thus \(\varphi(x_0) \leq x_0 \leq \varphi(x_0)\), so \(x_0 = \varphi(x_0)\) and the lemma would be proven. So we need only check that \(A\) is a chain, i.e., satisfies (iii). To do so let

\[
P = \{p \in A \mid \text{if } y \in A \text{ and } y < p, \text{ then } \varphi(y) \leq p\}.
\]

We have \(P \subseteq A\) and wish to show that \(A = P\). To show this, fix \(p \in P\) and define

\[
B(p) = \{z \in A \mid \text{either } z \leq p \text{ or } z \geq \varphi(p)\},
\]

i.e., \(B(p)\) contains no \(z \in A\) properly between \(\varphi(p)\) and \(p\).

**Step 1.** \(A = B(p)\):

We show that (i), (ii), (iii) hold for \(B(p)\). Let \(z \in A\).

(i): For all \(y \in A\), we have \(a \leq y\) and by (ii) for \(A\), we know that \(\varphi(y) \in A\). So \(a \in B(p)\).

(ii): If \(z < p\), then \(\varphi(z) \leq p\) by definition of \(P\). If \(z = p\), then \(\varphi(z) = \varphi(p) \geq p\). If \(z \geq \varphi(p)\), then \(\varphi(z) \geq z \geq \varphi(p) \geq p\). It follows that if \(z \in B(p)\), then \(\varphi(z)\) lies in \(B(p)\).

(iii): Let \(C\) be a chain in \(B(p)\) and \(c_0\) the least upper bound of \(C\). Then \(c_0\) lies in \(A\). Suppose that \(c \leq p\) for all \(c \in C\). Then \(c_0 \leq p\), hence \(\varphi(c_0) \leq p\) by (ii), so \(c_0\) lies in \(B(p)\). Otherwise, there exists a \(c \in C\) satisfying \(c \geq \varphi(p)\). Then \(c_0 \geq c \geq \varphi(p)\). This also implies that \(c_0\) lies in \(B(p)\) using (ii). It follows that the least upper bound of \(C\) lies in \(B(p)\).

We have shown that \(B(p) \in B_a\). Since \(B(p) \subseteq A\), we must have \(B(p)\). This establishes Step 1.

**Step 2.** \(P = A\).

We show that (i), (ii), (iii) hold for \(A\) and \(P\). Let \(z \in A\).

(i): We have \(a \in P\), since no \(y \in A\) can satisfy \(y < a\).

(ii): Let \(p \in P\). We want to show that \(\varphi(p) \in P\), i.e., if \(z \leq \varphi(p)\), then \(\varphi(z) \leq \varphi(p)\). We know that \(\varphi(p) \in A\) by (ii) for \(A\). Suppose that \(z < \varphi(p)\). Since \(A = B(p)\), we must have \(z \leq p\). If \(z < p\),
then by definition of $P$, we have $\varphi(z) \leq p \leq \varphi(p)$; and if $z = p$, then $\varphi(z) = \varphi(p) \leq \varphi(p)$. Hence $\varphi(p) \in P$ and $\varphi(P) \subseteq P$.

(iii): Let $C$ be a chain in $P$. Then the least upper bound $c_0$ of $C$ lies in $A$. Let $z \in A$. Suppose that we have $z < c_0$. If $z \in P$, then $c_0 \geq \varphi(z)$ by (ii) for $P$. If $z \notin P$ then for each $c \in C$, either $z < c$ or $z \geq \varphi(c) \geq c$ as $A = B(p)$. But $c_0$ is the least upper bound of $C$, so we cannot have $z \geq c$ for all $c \in C$. Therefore, $c_0 \in P$. It follows that $P \in B_a$. Since $P \subseteq A$, we must have $P = A$. Finally, if $x \in A$ and $y \in A$, then either $x \leq y$ or $x \geq \varphi(z) \geq z$. Therefore, $A$ is a chain. This proves the lemma. □

**Proposition A.11.** Suppose the Axiom of Choice holds. Then Zorn’s Lemma holds.

**Proof.** Let $S$ be strongly inductive and $S$ the collection of all non-empty subsets of $S$. Let $\psi : S \to S$ be the function satisfying $\psi(A) \in A$ for all $A \in S$. Define $\varphi : S \to S$ by $\varphi(x) =$

$$
\begin{cases}
  x & \text{if } x \text{ is a maximal element in } S \\
  \psi(\{z \mid \psi(\{z \mid z > x\})\}) & \text{if } x \text{ is not a maximal element in } S.
\end{cases}
$$

Then $\varphi(x) \leq x$ for all $x \in S$. Hence by Lemma A.10, there is an $x_0$ in $S$ such that $\varphi(x_0) = x_0$. Then $x_0$ is a maximal element of $S$. Therefore, every non-empty strongly inductive poset has a maximal element. By Lemma A.3, Zorn’s Lemma holds. □
APPENDIX B

Matrix Representations

We review the notion of matrix representation of a linear transformation relative to a pair of ordered bases and what happens when we change bases. We do this in somewhat more generality than in linear algebra. All proofs, however, are the same, hence omitted. Throughout $R$ will be a commutative ring. [So if $R$ is a field, and $R$-module is just a vector space over $R$ and an $R$-homomorphism is just a linear transformation.]

Setup.

Let $R$ be a commutative ring.
$V$ and $W$ are free $R$-modules of rank $n$ and $m$ respectively.

$\mathcal{B} := \{v_1, ..., v_n\}$ is an ordered basis for $V$.
$\mathcal{B}'$ is a second ordered basis of $V$.

$\mathcal{C} := \{w_1, ..., w_m\}$ is an ordered basis for $W$.
$\mathcal{C}'$ is a second ordered basis of $W$.

$T : V \to W$ is an $R$-homomorphism.

Let $v$ be an element in the free $R$-module $V$. Then there exist unique scalars $\alpha_1, ..., \alpha_n$ in $R$ such that

$$v = \alpha_1 v_1 + \cdots + \alpha_n v_n.$$ 

The scalars $\alpha_1, ..., \alpha_n$ are called the coordinates of $v$ relative to the (ordered) basis $\mathcal{B}$.

Let

$$[v]_{\mathcal{B}} = \begin{pmatrix} \alpha_1 \\ \vdots \\ \alpha_n \end{pmatrix}$$

the coordinate matrix of $v$ relative to the basis $\mathcal{B}$ and

$$V_{\mathcal{B}} := \{ [v]_{\mathcal{B}} \mid v \in V \} = R^{n \times 1}.$$ 

We have an $R$-isomorphism

$$V \to V_{\mathcal{B}} \text{ given by } v \mapsto [v]_{\mathcal{B}}.$$
Similarly, if \( w \) is a vector in the free \( R \)-module \( W \), then there exist unique scalars, \( \beta_1, \ldots, \beta_m \), in \( R \) such that
\[
w = \beta_1 w_1 + \cdots + \beta_m w_m.
\]
Let
\[
[w]_C = \begin{pmatrix} \beta_1 \\ \vdots \\ \beta_m \end{pmatrix}
\]
the coordinate matrix of \( w \) relative to the basis \( C \) and \( W_C := \{ [w]_C \mid w \in W \} = R^{m \times 1} \).

Examples B.1. 1.
\[
[v_1]_B = \begin{pmatrix} 1 \\ 0 \\ \vdots \\ 0 \end{pmatrix}, \ldots, [v_n]_B = \begin{pmatrix} 0 \\ \vdots \\ 0 \\ 1 \end{pmatrix}.
\]
2. Suppose that \( n = 3 \). Then
\[
[2v_1 - 3v_3]_B = \begin{pmatrix} 2 \\ 0 \\ -3 \end{pmatrix}, \quad [7v_2 - v_3]_B = \begin{pmatrix} 0 \\ 7 \\ -1 \end{pmatrix}.
\]
We turn to the \( R \)-homomorphism \( T : V \to W \). By the Universal Property of Free Modules, there exists a unique matrix
\[
[T]_{B,C} \in R^{m \times n}
\]
called the matrix representation of \( T \) relative to the ordered bases \( B, C \) that satisfies
\[
[T]_{B,C}[v]_B = [T(v)]_C \quad \forall v \in V.
\]
If \( V = W \) and \( B = C \), we let \( [T]_B = [T]_{B,C} \).

The matrix \( [T]_{B,C}[v]_B = [T(v)]_C \) is computed as follows. Write \(Tv_j\) in the \( C \) basis. Then the coordinate matrix of \(Tv_j\) relative to the \( C \) basis is the \( j \)th column of the matrix \( [T]_{B,C} \), i.e., if
\[
Tv_j = \beta_{1j}w_1 + \cdots + \beta_{mj}w_m \quad \text{then} \quad [T(v_j)]_C = \begin{pmatrix} \beta_{1j} \\ \vdots \\ \beta_{mj} \end{pmatrix}
\]
and this is the \( j \)th column of \( [T]_{B,C} \). It is convenient’ to write
Definition B.2. A matrix of the form $[1_V]_{B,B'}$ is called a change of basis matrix.

It arises by writing the elements in the $B$ basis in terms of the elements in the $B'$ basis.

The main results are:

**Proposition B.3.** Let $V$ be an $R$-free module of rank $n$. Let $B$ and $B'$ be ordered bases for $V$. Then $[1_V]_{B,B'}$ is an invertible matrix and

$$[1_V]_{B,B'}^{-1} = [1_V]_{B',B}$$

**Remarks B.4.** It can be shown that

$$\text{GL}_n R = \{[T]_{B,C} | T \text{ an isomorphism, } B, C \text{ bases for } V\}.$$ 

Since $R$ is commutative, determinants exist. It can also be shown that

$$\text{GL}_n R = \{A \in M_n R | \det(A) \in R^\times\}$$

**Theorem B.5.** Let $V$, $W$, and $X$ be finitely generated free $R$-modules. Let $B$, $C$, and $D$ be ordered bases for $V$, $W$, and $X$, respectively. Let $T: V \rightarrow W$ and $S: W \rightarrow X$ be $R$-homomorphisms. Then

$$[S \circ T]_{B,D} = [S]_{C,D}[T]_{B,C}$$

where the right hand side is matrix multiplication.

**Theorem B.6.** Let $V$ be a finitely generated free $R$-module and $B$ an ordered basis for $V$. Then

$$\text{End}_R V \rightarrow M_n R \text{ via } T \mapsto [T]_B$$

is a ring isomorphism and induces a group isomorphism

$$\text{Aut}_R V \rightarrow \text{GL}_n R.$$ 

**Theorem B.7.** (Change of Basis Theorem.) Let $V$ and $W$ be finitely generated $R$-free modules. Let $B$ and $B'$ be ordered bases for $V$ and $C$ and $C'$ be ordered basis for $W$. Let $T: V \rightarrow W$ be an $R$-homomorphism. Then


The Change of Basis Theorem states that the following diagram commutes
Note that the inverses of the change of bases matrices go in the reverse direction. One can fill in more of the diagram. For example, the maps from the diagonals can also be read off, e.g.,

\[
\begin{align*}
[T]_{B,B'} &= [T]_{B',B'} [1_V]_{B,B'} \\
[T]_{B',C'} &= [1_V]_{B',C'} [T]_{B,B'} [1_V]_{B,B'}^{-1}
\end{align*}
\]

**Warning B.8.** Usually \( T \) is not an isomorphism, so \([T]_{B,C'}\) is not invertible. So you cannot reverse arrows having \( T \) in them. If \( T \) is an isomorphism, then the matrix representation of \( T^{-1} \) is the inverse of the corresponding matrix representation of \( T \).

Historically, one defined equivalence relationships between matrices. We define these and tell how they are related to \( R \)-homomorphisms of finitely generated free \( R \)-modules.

**Definition B.9.** If \( A, B \in M_n R \), we say that they are *similar* and write \( A \sim B \) if there is an invertible matrix \( C \in GL_n R \) such that \( A = C^{-1} B C \).

Clearly, \( \sim \) is an equivalence relation. An important problem is to find good representatives for the classes under this equivalence relation when \( R \) is a nice ring, e.g., a PID. This leads to the study of Rational Canonical Forms and Jordan Canonical Forms in linear algebra.

**Theorem B.10.** Let \( A, B \in M_n R \). Then \( A \) is similar to \( B \) in \( M_n R \) if and only if there exist a free \( R \)-module \( V \) of rank \( n \) with bases \( B, B' \) and \( T \in \text{End}_R V \) such that \( A = [T]_B \) and \( B = [T]_{B'} \).

**Definition B.11.** If \( A, B \in R^{m \times n} \), we say that they are *equivalent* and write \( A \simeq B \) if there is an invertible matrices \( P \in GL_m R \) and \( Q \in GL_n R \) such that \( A = PBQ \).

Clearly, \( \simeq \) is an equivalence relation. An important problem is to find good representatives for the classes under this equivalence relation when \( R \) is a nice ring, e.g., a PID. This leads to the study of Smith Normal Forms in linear algebra.

**Theorem B.12.** Let \( A, B \in R^{m \times n} \). Then \( A \simeq B \) in \( R^{m \times n} \) if and only if there exist free \( R \)-modules \( V, W \) of rank \( n, m \), respectively, with bases
$\mathcal{B}, \mathcal{B}'$ for $V$ and bases $\mathcal{C}, \mathcal{C}'$ for $W$ and an $R$-homomorphism $T : V \to W$ such that $A = [T]_{\mathcal{B}, \mathcal{C}}$ and $B = [T]_{\mathcal{B}', \mathcal{C}'}$. 
APPENDIX C

Smith Normal Form over a Euclidean Ring

Let $R$ be a domain and $A \in R^{m \times n}$. Recall that $A$ is in Smith Normal Form if $A$ is a matrix of the form

$$
\begin{pmatrix}
q_1 & 0 & \ldots \\
0 & q_2 & \\
& \ddots & \\
& & q_r & 0 \\
& & & \ddots \\
0 & & & & \ddots \\
\end{pmatrix}
$$

with $q_1 \mid q_2 \mid q_3 \mid \cdots \mid q_r$ in $R$ and $q_r \neq 0$, i.e., the diagonal entries of $A$ are $q_1, \ldots, q_r, 0, \ldots, 0$ with $q_1 \mid q_2 \mid q_3 \mid \cdots \mid q_r$ in $R$ and $q_r \neq 0$ and all entries off the diagonal are 0.

Let $R$ be a domain and $A \in M_n(R)$. We say that $A = (a_{ij})$ is an elementary matrix if

(i) **Type I:** if there exists $\lambda \in R$ and $l \neq k$ such that

$$
a_{ij} = \begin{cases} 1 & \text{if } i = j \\ \lambda & \text{if } (i, j) = (k, l) \\ 0 & \text{otherwise} \end{cases}
$$

(ii) **Type II:** if there exists $k \neq l$ such that

$$
a_{ij} = \begin{cases} 1 & \text{if } i = j \neq l \text{ and } i = j \neq k \\ 0 & \text{if } i = j = l \text{ or } i = j = k \\ 1 & \text{if } (k, l) = (i, j) \text{ or } (k, l) = (j, i) \\ 0 & \text{otherwise} \end{cases}
$$

(iii) **Type III:** if there exists a unit $u$ in $R$ and $l$ such that
\[ a_{ij} = \begin{cases} 
1 & \text{if } i = j \neq l \\
u & \text{if } i = j = l \\
0 & \text{otherwise} 
\end{cases} \]

**Remarks C.1.**

1. Let \( A \in \mathbb{R}^{m \times n} \). Multiplying \( A \) on the left (respectively right) by a suitable size

   (i) Type I is adding a multiple of a row (respectively column) of \( A \) to another row (respectively column) of \( A \).

   (ii) Type II is interchanging two rows (respectively columns) of \( A \).

   (iii) Type III is multiplying a row (respectively column) of \( A \) by a unit.

2. All elementary matrices are invertible.

Recall if \( R \) is a ring, two \( m \times n \) matrices \( A \) and \( B \) in \( R^{m \times n} \) are called equivalent if there exist invertible matrices \( P \in \text{GL}_m(R) \) and \( Q \in \text{GL}_n(R) \) such that \( A = PAQ \).

**Theorem C.2.** Let \( R \) be a euclidean ring, \( A = (a_{ij}) \in \mathbb{R}^{m \times n} \). Then \( A \) is equivalent to a matrix in Smith Normal Form. Moreover, there exist matrices \( P \in \text{GL}_m(R) \) and \( Q \in \text{GL}_n(R) \), each a product of matrices of Type I, Type II, and Type III, such that \( PAQ \) is in Smith Normal Form.

**Proof.** The proof will, in fact, be an algorithm to find a Smith Normal Form for \( A \). Let \( \delta \) be the euclidean function on \( R \). If \( A = 0 \) there is nothing to do, so assume that \( A \neq 0 \).

**Step 1.** Choose \( a = a_{ij} \neq 0 \) such that \( \delta(a) \) is minimal among all the \( \delta(a_{lk}), a_{lk} \neq 0 \). Put \( a \) in the \((1, 1)\) spot using matrices of Type II. In particular, we may assume that \( a = a_{11} \).

[If \( a \) is a unit in \( R \), use a Type III matrix to make \( a = 1 \).]

**Step 2.** If \( a \mid a_{ij} \) for all \( i \) and \( j \), use Type I matrices to transform \( A \) into a matrix of the form

\[
\begin{pmatrix}
a & 0 & \ldots & 0 \\
0 & \ddots & & \\
0 & & \ddots & A_1 \\
\end{pmatrix}
\]

Note that \( a \) divides every entry of \( A_1 \) [Check]. If a non-zero entry of \( A_1 \) has smaller \( \delta \) value, \( \delta(a) \) is not minimal so go back to Step 1.

[As \( \delta a \) is a non-negative integer this cannot happen infinitely often.]
If $\delta a$ is still minimal, take $A_1$ and go back to Step 1.

[Note. If this occurs, by induction there exist matrices $Q_1$, $P_1$ such that $P_1A_1Q_1$ is in Smith Normal Form. Let

\[
P = \begin{pmatrix} 1 & 0 & \ldots & 0 \\ 0 & \ddots & \ddots & \vdots \\ \vdots & & P_1 & \\ 0 & & & \end{pmatrix} \quad \text{and} \quad Q = \begin{pmatrix} 1 & 0 & \ldots & 0 \\ 0 & \ddots & \ddots & \vdots \\ \vdots & & Q_1 & \\ 0 & & & \end{pmatrix}
\]

then

\[
P \begin{pmatrix} a & 0 & \ldots & 0 \\ 0 & \ddots & \ddots & \vdots \\ \vdots & & A_1 & \\ 0 & & & \end{pmatrix} Q
\]

is in Smith Normal Form.]

**Step 3.** Step 2 does not apply and there exists an entry $b = a_{ij}$ in either the first row or first column such that $a \nmid b$:

Write $b = qa + r$ in $R$ with $r \neq 0$ and $\delta(r) < \delta(a)$. Use Type I matrices to change $A$ into a matrix with $r$ in it. Since $\delta(a)$ is not longer minimal, go back to Step 1.

[Since $\delta(a)$ is a non-negative integer and $\delta(r) < \delta(a)$, this must eventually stop.]

**Step 4.** Neither Step 2 nor Step 3 apply. Thus $a \mid a_{ij}$ whenever $i = 1$ or $j = 1$:

Use Type I matrices to convert $A$ to

\[
\begin{pmatrix} a & 0 & \ldots & 0 \\ 0 & \ddots & \ddots & \vdots \\ \vdots & & A_1 & \\ 0 & & & \end{pmatrix}
\]

Now one of the following occurs.

1. There exists a non-zero entry $b$ in $A_1$ such that $\delta(b) < \delta(a)$. So $\delta(a)$ is no longer minimal. Go back to Step 1.

2. $a \mid b$ for all entries $b$ in $A_1$. This is impossible — You should have been in Step 2.

   [No matter, take $A_1$ and go to Step 1.]

3. There exists an entry $b$ in $A_1$ such that $a \nmid b$:

   Write $b = qa + r$ in $R$ with $r \neq 0$ and $\delta(r) < \delta(a)$. Use Type I matrices to get $b$ into the first column. (This does not change
the (1, 1) entry $a$.) Now use Type I matrices to change $b$ to $r$. Since $\delta(a)$ is no longer minimal, go back to Step 1.

Clearly this algorithm yields a Smith Normal Form of $A$.

\[ \square \]

Remarks C.3. 1. If $R$ is a ring, two $m \times n$ matrices $A$ and $B$ in $R^{m \times n}$ are called equivalent if there exist invertible matrices $P \in \text{GL}_m(R)$ and $Q \in \text{GL}_n(R)$ such that $A = PAQ$. Compare this to change of basis in linear algebra of matrix representations of linear transformations. So the theorem says if $R$ is a euclidean ring then any $m \times n$ matrix over $R$ is equivalent to a matrix in Smith Normal Form.

2. Note that we do not really need Type III matrices in the above. We only used them to transform those diagonal entries $q_i$ that were units into one.

3. We did not really need Step 2 as it is incorporated into Step 4, but it is useful to isolate it to complete the induction step of the proof.

As stated in Theorem 42.8 and proved, a Smith Normal Form of a matrix over a PID is unique up to units. We shall assume it here. Its proof depends only on an elementary determinant argument and could easily be done here, but we shall not repeat it in this appendix.

In Section 44, Smith Normal Forms are used in understanding parts of linear algebra. In this appendix, we shall prove (most) of Theorem 44.11 without the development needed in Section 44 used to prove it there. Rather we shall only use the Smith Normal Form of a matrix over the euclidean domain $F[t]$ with $F$ a field and a generalization of the division algorithm for $F[t]$.

Let $V$ be a finite dimensional vector space over a field $F$ and $B$ an ordered basis for $V$. Classifying linear operators $T : V \rightarrow V$ is equivalent to determining $[T]_B$ up to matrix similarity. So we turn to the matrix formulation.

Let $A_0$ be a matrix in $\mathbb{M}_n(F)$ and $A = tI - A_0$, a matrix in $\mathbb{M}_n(F[t])$. This matrix is called the characteristic matrix of $A_0$. As $\det A$ is the characteristic polynomial of $A_0$, hence nonzero, there is a Smith Normal Form of $A$ of the form $\text{diag}(1, \ldots, 1, q_1, \ldots, q_r)$ for some nonconstant polynomials $q_1, \ldots, q_r$ in $F[t]$ and some unique $r$. As the Smith Normal Form of $A$ is unique up to units (by Theorem 42.8), using Type III operations, we may assume that the polynomials $q_1, \ldots, q_r$ are all monic. This is will be called the Smith Normal Form of $A$ and the unique monic polynomials $q_1, \ldots, q_r$ will be called the invariant factors of $A$. If $B_0$ is another matrix in $\mathbb{M}_n(F)$ and $B = tI - B_0$, we want to
show that $A$ is equivalent to $B$ in $\mathbb{M}_n(F[t])$ if and only if $A_0$ is similar to $B_0$. As mentioned above this is also shown in Section 44, but our proof is not only more elementary but interesting in its own right.

We need a special case of the division algebra applied to matrices in $\mathbb{M}_n(F[t])$. We know that the division algorithm says if $g \in F[t] \setminus \{0\}$ and $a \in F$, then $g = (t-a)h + r$ for some $h$ in $F[t]$ and $r \in F$ — indeed $r = g(a)$ by the Remainder Theorem. We want a similar result for matrices in $\mathbb{M}_n(F[t])$. However, dividing on the right and dividing on the left makes a difference, i.e., we may have $G = AQ_1 + R_1 = Q_2 A + R_2$, but $Q_1 \neq Q_2$ and $R_1 \neq R_2$ in the analogue of the above. However, the same algorithm in the commutative case can be mimicked as we shall now show.

We use the following notation:

**Notation C.4.** Let $F$ be a field, $A_i \in \mathbb{M}_n(F)$, $i = 0, \ldots, n$. The polynomial matrix

$$A_n t^n + A_{n-1} t^{n-1} + \cdots + A_0$$

will denote the matrix

$$A_n(t^n I) + A_{n-1}(t^{n-1} I) + \cdots + A_0$$

in $\mathbb{M}_n(F[t])$ with $t$ commuting with all elements of $F$ and $I$ the identity matrix. So if $A = (\alpha_{ij})$, then $At^n = (\alpha_{ij} t^n)$. Two polynomials matrices are the same if and only if their coefficients are equal, i.e., we are identifying $\mathbb{M}_n(F[t])$ and $(\mathbb{M}_n(F))[t]$.

**Lemma C.5.** Let $F$ be a field and $A_0$ an element in $\mathbb{M}_n(F)$. Set $A = tI - A_0$ in $\mathbb{M}_n(F[t])$. Suppose that $P = P(t)$ is an element in $\mathbb{M}_n(F[t])$. Then there exist matrices $M$ and $N$ in $\mathbb{M}_n(F[t])$ and matrices $R$ and $S$ in $\mathbb{M}_n(F)$ satisfying

1. $P = AM + R$.
2. $P = NA + Q$.

**Proof.** We show (1), the second being similar. Let

$$m = \max_{l,k}(\deg P_{lk}) \quad \text{and} \quad P_{ij} = \alpha_{ij} t^m + \text{lower terms in } F[t]$$

for all appropriate $l,k,i,j$. In particular, we have

$$\alpha_{ij} = \begin{cases} \text{lead} P_{ij} & \text{if } \deg P_{ij} = m \\ 0 & \text{if } \deg P_{ij} < m. \end{cases}$$

Let $(\alpha_{ij}) \in \mathbb{M}_n(F)$ and define a polynomial matrix

$$P_{m-1} = (\alpha_{ij}) t^{m-1} = (\alpha_{ij} t^{m-1}) = (t^{m-1} \alpha_{ij}).$$
Every entry in
\[ AP_{m-1} = (tI - A_0)(\alpha_{ij})t^{m-1} = (\alpha_{ij})t^m - A_0(\alpha_{ij})t^{m-1} \]
has degree at most \( m \) and the \( t^m \)-coefficient of \( (AP_{m-1})_{ij} \) is \( \alpha_{ij} \).

Therefore, \( P - AP_{m-1} \) has polynomial entries of degree at most \( m - 1 \).

Applying the same argument to \( P - AP_{m-1} \) produces a matrix \( P_{m-2} \)
in \( \mathbb{M}_n(F[t]) \) with all the polynomial entries of \( (P - AP_{m-1}) - AP_{m-2} \) having degree at most \( m - 2 \). Continuing in this way, we construct matrices \( P_{m-3}, \ldots, P_0 \) such that if \( M := P_{m-1} + P_{m-2} + \cdots + P_0 \), then \( R := P - AM \) has only constant entries, i.e., \( R \in \mathbb{M}_n(F) \). Therefore, \( P = AM + R \) as needed.

We use the lemma to prove our theorem:

**Theorem C.6.** Let \( F \) be a field and \( A_0, B_0 \) be elements in \( \mathbb{M}_n(F) \). Set \( A = tI - A_0 \) and \( B = tI - B_0 \). Then \( A \) and \( B \) are equivalent in \( \mathbb{M}_n(F[t]) \) if and only if \( A_0 \) and \( B_0 \) are similar in \( \mathbb{M}_n(F) \).

**Proof.** \((\Leftarrow)\) : If \( A_0 = P B_0 P^{-1} \) for some \( P \in \text{GL}_n(F) \), then
\[ P(tI - A_0)P^{-1} = PtIP^{-1} - PA_0P^{-1} = tI - B_0 = B. \]
So \( B = PAP^{-1} \) and \( B \) is equivalent to \( A \).

\((\Rightarrow)\) : Suppose that there exist matrices \( P_1 \) and \( Q_1 \) in \( \text{GL}_n(F[t]) \), (each is a product of elementary matrices by Theorem C.2) satisfying
\[ B = tI - B_0 = P_1AQ_1 = P_1(tI - A_0)Q_1. \]

By the Lemma C.5, we can write
\( (i) \) \( P_1 = BP_2 + R \) with \( P_2 \in \mathbb{M}_n(F[t]) \) and \( R \in \mathbb{M}_n(F) \).
\( (ii) \) \( Q_1 = Q_2B + S \) with \( Q_2 \in \mathbb{M}_n(F[t]) \) and \( S \in \mathbb{M}_n(F) \).

As \( B = P_1AQ_1 \) with \( P_1, Q_1 \in \text{GL}_n(F[t]) \), we also have
\( (iii) \) \( P_1A = BQ^{-1}_1. \)
\( (iv) \) \( AQ_1 = P^{-1}_1B. \)

Therefore, we have
\[ B = P_1AQ_1 = (BP_2 + R)AQ_1 = BP_2AQ_1 + RAQ_1 \]
\[ = BP_2P_1^{-1}B + RAQ_1 = BP_2P_1^{-1}B + RA(Q_2B + S) \]
\[ = BP_2P_1^{-1}B + RAQ_2B + RAS, \]
i.e., we have
\( (v) \) \( B = BP_2P_1^{-1}B + RAQ_2B + RAS. \)
We work on the middle term $RAQ_2B$ of this equation. By (i), we have

$$R = P_1 - BP_2,$$

so plugging this in for $RAQ_2B$, yields

$$RAQ_2B = (P_1 - BP_2)AQ_2B = P_1AQ_2B - BP_2AQ_2B \tag{iii}$$

$$= BQ_1^{-1}Q_2B - BP_2AQ_2B = B[Q_1^{-1}Q_2 - P_2AQ_2]B,$$

i.e., we have

$$(vi) \ RAQ_2B = B[Q_1^{-1}Q_2 - P_2AQ_2]B.$$ 

Plugging $(vi)$ into $(v)$ yields

$$B = BP_2P_1^{-1}B + RAQ_2B + RAS \tag{v}$$

$$= BP_2P_1^{-1}B + B[Q_1^{-1}Q_2 - P_2AQ_2]B + RAS \tag{vi}$$

$$= B[P_2P_1^{-1} + Q_1^{-1}Q_2 - P_2AQ_2]B + RAS.$$ 

Let

$$T = P_2P_1^{-1} + Q_1^{-1}Q_2 - P_2AQ_2.$$ 

So we have

$$(vii) \ B = BTB + RAS.$$ 

Next we look at the degree of the polynomial entries in these three matrices in $(vii)$. Certainly,

$$(viii) \text{ Every entry of } B = tI - B_0 \text{ has degree at most one and every entry of } RAS = R(tI - A_0) \text{ has degree at most one.}$$ 

So we must look at $BTB$. Write

$$T = T_m t^m + T_{m-1} t^{m-1} + \cdots + T_0 \text{ with } T_0, \ldots, T_m \in M_n(F).$$ 

Then

$$BTB = (tI - B_0)(T_m t^m + \cdots + T_0)(tI - B_0)$$

$$= T_m t^{m+2} + \text{ lower terms in } t.$$ 

Comparing coefficients of the matrix polynomials in $(vii)$ and using $(viii)$ shows that $T_m = 0$. It follows that $T = 0$. Thus $(vii)$ becomes

$$(C.7) \ tI - B_0 = B = BTB + RAS = RAS = R(tI - A_0)S$$

$$= RSt + RA_0S.$$
Comparing coefficients of the polynomial matrices in Equation C.7 now implies that $I = RS$, i.e., $S = R^{-1}$, and $B_0 = RA_0S$. In particular,

$$B_0 = RA_0S = RA_0R^{-1}.$$ 

Hence $B_0$ is similar to $A_0$ as needed. \hfill \Box

We therefore have, as Smith Normal Form is unique,

**Theorem C.8.** Let $F$ be a field and $A_0, B_0$ be elements in $\mathbb{M}_n(F)$. Set $A = tI - A_0$ and $B = tI - B_0$. Then the following are equivalent.

1. $A$ and $B$ are equivalent in $\mathbb{M}_n(F[t])$.
2. $A_0$ and $B_0$ are similar in $\mathbb{M}_n(F)$.
3. $A$ and $B$ have the same invariant factors.

**Corollary C.9.** Let $F$ be a field and $A$ a matrix in $\mathbb{M}_n(F)$. Then $A$ is similar to $A^t$.

**Proof.** They have the same Smith Normal Form. \hfill \Box

**Exercise C.10.** Let $R$ be a commutative ring. Let $E_n(R)$ be the subgroup of $\text{GL}_n(R)$ generated by all matrices of the form $I + \lambda$ where $\lambda$ is a matrix with precisely one non zero entry and this entry does not occur on the diagonal. Suppose that $R$ is a euclidean ring. Show that $\text{SL}_n(R) = E_n(R)$. 

APPENDIX D

Symmetric Bilinear Forms

In this appendix, we establish Sylvester’s theorem about the signature of a symmetric matrix over the real numbers. We shall need its interpretation over a subfield $F$ of the reals using the usual ordering $<$ of $\mathbb{R}$ restricted to $F$, i.e., we have a disjoint union $F = F^+ \cup \{0\} \cup F^-$ with $F^+ := \{x \in F \mid x > 0\}$ and $F^- := \{x \in F \mid x < 0\}$. The proof of Sylvester’s Theorem will then hold in the more general situation needed in Section 73.

We begin with a short discussion of symmetric bilinear spaces. For convenience, we restrict our attention to finite dimensional symmetric bilinear forms over a field of characteristic different from two.

**Definition D.1.** Let $V$ be a finite dimensional vector space over a field $F$ of characteristic different from two and $B : V \times V \to F$ a symmetric bilinear form, i.e., linear in each variable and $B(x, y) = B(y, x)$ for all $x, y \in V$. We call $(V, B)$ a symmetric bilinear space. If $W \subset V$ is a subspace of $V$, then $B|_{W \times W} : W \times W \to F$ is a symmetric bilinear form, and we call $(W, B|_{W \times W})$ a (bilinear) subspace of $(V, B)$. We say $x, y$ in $V$ are orthogonal (rel $B$) if $B(x, y) = 0$. We say two subspaces $W_1, W_2$ of $V$ are orthogonal if $B(x, y) = 0$ for all $x \in W_1$, $y \in W_2$ and, if this is the case, we write $W_1 \perp W_2$. Note if this is the case, then $W_1 + W_2 = W_1 \oplus W_2$ if and only if $W_1 \cap W_2 = \{0\}$ (by dimension count). If $V = W_1 \oplus W_2$ and $W_1 \perp W_2$, then we write $(V, B) = (W_1, B|_{W_1 \times W_1}) \perp (W_2, B|_{W_2 \times W_2})$, and call it an orthogonal decomposition of $(V, B)$.

**Lemma D.2.** Let $F$ be a field of characteristic different from two and $(V, B)$ a finite dimensional symmetric bilinear space. If $W$ is a subspace of $V$, then

$$W^\perp := \{v \in V \mid B(v, w) = 0 \text{ for all } w \in W\}$$

is a subspace of $V$ and $(W^\perp, B|_{W^\perp \times W^\perp})$ is a subspace of $(V, B)$. In particular, $(V, B) = (W, B|_{W \times W}) \perp (W^\perp, B|_{W^\perp \times W^\perp})$ if and only if $W \oplus W^\perp = V$.

**Proof.** If $v_1, v_2 \in W^\perp$ and $\alpha \in F$, then

$$B(\alpha v_1 + v_2, w) = \alpha B(v_1, w) + B(v_2, w) = 0$$
for all \( w \in W \), i.e., \( W^\perp \) is a subspace. The second statement is clear. \( \square \)

**Remark D.3.** Let \( F \) be a field of characteristic different from two and \((V, B)\) a finite dimensional symmetric bilinear space. Then
\[
B(x, y) = \frac{1}{2} [(B(x + y, x + y) − B(x, x) − B(y, y)].
\]
Since \( 2 \neq 0 \) under the characteristic assumption on \( F \), it follows that \( B \) is determined by the composition \( q_B := B \circ \Delta : V \to V \), where \( \Delta : V \to V \times V \) given by \( \Delta(v) = (v, v) \) for all \( v \in V \) is the diagonal map. The map \( q_B \) is called the associated quadratic form of \( B \). In particular, if \( B(x, y) \neq 0 \), then by Remark D.3, \( B(z, z) \neq 0 \) for \( z \) at least one of \( x, y, x + y \) since \( 2 \neq 0 \). This is the major reason that we are assuming that the characteristic of \( F \) is not two.

**Definition D.4.** Let \((V, B)\) be finite dimensional symmetric bilinear space. Let \( \text{rad}(B) := \{ v \mid B(v, w) = 0 \text{ for all } w \in V \} \), called the radical of \((V, B)\). So \( \text{rad}(B) = V^\perp \). If \( \text{rad}(B) = 0 \), we say that \((V, B)\) is regular.

Let \((V, B)\) be a finite dimensional symmetric bilinear space over \( F \), a field of characteristic different from two. Then for every \( v \in V \), the map \( B_v : V \to F \) defined by \( x \mapsto B(v, x) \) is a linear functional, so defines a map \( \widehat{B} : V \to V^* \) by \( v \mapsto B_v \), with \( V^* \) the dual space of \( V \). This map is checked to be an isomorphism if and only if \( V \) is regular.

**Definition D.5.** Let \( F \) be a field of characteristic different from two and \((V, B), (V', B')\) finite dimensional symmetric bilinear spaces over \( F \). A vector space isomorphism \( \varphi : V \to V' \) is called an isometry if \( B'(\varphi(x), \varphi(y)) = B(x, y) \) for all \( x, y \in V \). If \( \varphi \) is an isometry we write \((V, B) \cong (V', B')\).

Let \( F \) be a field of characteristic different from two and \((V, B)\) a finite dimensional symmetric bilinear space. Let \( \overline{\varphi} : V \to V / \text{rad}(B) \) be the canonical epimorphism. Define \( \overline{B} \) to be the bilinear form on \( \overline{V} \) determined by \( \overline{B}((v_1, v_2)) := B(v_1, v_2) \) for all \( v_1, v_2 \in V \). Then \( \overline{(V, B)} \) is a regular symmetric bilinear space, called the induced bilinear space.

**Lemma D.6.** Let \( F \) be a field of characteristic different from two and \((V, B)\) a finite dimensional symmetric bilinear space. Let \( W \) be any subspace of \( V \) satisfying \( V = \text{rad}(B) \oplus W \). Then \((W, B|_{W \times W})\) is regular and
\[
(V, B) = (\text{rad}(B), B|_{\text{rad}(B) \times \text{rad}(B)}) \perp (W, B|_{W \times W})
\]
with \( (V, B|_{W \times W}) \cong (\overline{V}, \overline{B}) \), the form induced on \( V / \text{rad}(B) \). Moreover, \((W, B|_{W})\) is unique up to isometry.
Proof. The first statement follows, as \( \cdot : V \to V \) induces a vector space isomorphism \( V/W \to V \). Since any isometry of symmetric bilinear spaces takes radicals to radicals, the uniqueness statement follows. \( \Box \)

We leave the following lemma as an exercise.

**Lemma D.7.** Let \( F \) be a field of characteristic different from two, \((V,B)\) a finite dimensional symmetric bilinear over \( F \), and \( W \subset V \) a subspace. Define \( f_W : V \to W^* \) by \( v \mapsto (B|_W)v \). Then \( W^\perp = \ker f_W \), i.e., we have an an exact sequence of vector spaces

\[
0 \to W^\perp \to V \xrightarrow{f_W} W^*,
\]

where \( \text{inc} \) is the inclusion map. In particular, \( \dim W^\perp \geq \dim V - \dim W \).

Using the lemma, it is easy to determine when this is an equality.

**Proposition D.8.** Let \( F \) be a field of characteristic different from two and \((V,B)\) be a finite dimensional symmetric bilinear form. Let \( W \) be any subspace of \( V \). Then the following are equivalent:

1. \( W \cap \text{rad}(B) = 0 \).
2. \( f_W : V \to W^* \) is surjective.
3. \( \dim W^\perp = \dim V - \dim W \).

Proof. By linear algebra if \( T : V \to W \) is a linear transformation, there is a dual map of dual spaces \( T^\dagger : W^* \to V^* \) where \( T^\dagger \), called the transpose of \( T \) is the linear transformation defined by \( T^\dagger(g) = g \circ T \). Moreover, applied to compositions, it takes exact sequences to exact sequences (with the maps reversed upon taking the dual sequence). As \( V \) is finite dimensional, \( V = V^{**} \). It follows that (1) holds if and only if the transpose map \( (f_W)^\dagger : W \to V^* \) is injective if and only if the map \( f_W : V \to W^* \) is surjective (why?) if and only if (3) holds. \( \Box \)

A key observation is:

**Proposition D.9.** Let \( F \) be a field of characteristic different from two and \((V,B)\) a finite dimensional symmetric bilinear space. Suppose that \( W \) is a subspace of \( V \) with \((W,B|_{W \times W})\) regular. Then \((V,B) = (W,B|_{W \times W}) \perp (W^\perp,B|_{W^\perp \times W^\perp})\). In particular, if \((V,B)\) is also regular, then so is \((W^\perp,B|_{W^\perp \times W^\perp})\). \( \Box \)

Proof. By Proposition D.8, \( \dim W^\perp = \dim V - \dim W \), hence \( V = W \oplus W^\perp \). The result follows.
Corollary D.10. Let $F$ be a field of characteristic different from two and $(V, B)$ a finite dimensional symmetric bilinear space. If $v \in V$ satisfies $B(v, v) \neq 0$, then $(V, B) = (Fv, B|_{Fv \times Fv}) \perp (Fv^\perp, B|_{(Fv)^\perp \times (Fv)^\perp})$.

Theorem D.11. Let $F$ be a field of characteristic different from two and $(V, B)$ a finite dimensional symmetric bilinear space. Then $V$ has an orthogonal basis.

Proof. Let $\dim V = n$. If $V = \text{rad} V$ any basis works, so we may assume this is not so. In particular, there exists a vector $v \in V$ such that $B(v, v) \neq 0$ by Remark D.3. The result follows by induction on $\dim V$ in view of Corollary D.10. \square

If we had not put the condition that the field is of characteristic not two, then an orthogonal basis may not exist for a finite dimensional symmetric bilinear space.

As with linear operators on a finite dimensional vector space, we have matrix representations of finite dimensional symmetric bilinear spaces $(V, B)$. Let $B = \{v_1, \ldots, v_n\}$ be an ordered basis for $V$. Then the symmetric matrix $(B(v_i, v_j)) \in M_n(F)$ is called the matrix representation of $(V, B)$ relative to the basis $V$ and we denote it by $[B]_B$. We have $B(v, w) = v^t B w$ for all $v, w \in V$ (writing $w$ as a column matrix in the matrix equation) where $v^t$ is the transpose of $v$. In particular, $B$ an orthogonal basis if and only if $[B]_B$ is a diagonal matrix.

If $C$ is another ordered basis for $V$, computation shows that

$$[1_V]_B^t [B]_B [1_V]_B^C = [B]_C.$$

Let $F$ be a field of characteristic different from two and $(V, B)$ and $(V', B')$ two finite dimensional symmetric bilinear spaces with $B$, and $B'$ ordered bases for $V$ and $V'$, respectively. Then it is easy to see that $(V, B) \cong (V', B')$ if and only if $A^t[B]_B A = [B']_{B'}$ where $A$ is the matrix representation of the isomorphism $T : V \to V'$ taking the basis $B$ to $B'$.

We turn to the special case of interest in this book. Let $F$ be a subfield of $\mathbb{R}$, $V$ a finite dimensional vector space of dimension $n$ and $B : V \times V \to F$ a symmetric bilinear form. Let $<$ be the ordering on $F$ induced by $<$ on $\mathbb{R}$, so if $P := \{x \in F \mid x > 0\}$ and $-P := \{x \in F \mid x < 0\}$, we have $S = P \cup \{0\} \cup -P$ is a disjoint union. For example, if $F = \mathbb{R}$, then $P = \{x^2 \in F^\times \mid x \in F\}$.

Proposition D.12. Let $F$ be a subfield of $\mathbb{R}$, $V$ a finite dimensional vector space of dimension $n$ and $B : V \times V \to F$ a symmetric bilinear
form. Let $<$ be the ordering on $F$ induced by $<$ on $\mathbb{R}$. Then there exists an orthogonal decomposition $V = V_0 \oplus V_+ \oplus V_-$. with

\begin{align*}
B(v,v) &= 0 \quad \text{for all } v \in V_0 \\
B(v,v) &> 0 \quad \text{for all } 0 \neq v \in V_+ \\
B(v,v) &< 0 \quad \text{for all } 0 \neq v \in V_-.
\end{align*}

(D.13)

Moreover, the dimensions $\dim V_0$, $\dim V_+$, and $\dim V_-$ are independent of the orthogonal decomposition satisfying (D.13).

**Proof.** We first show that there exists an orthogonal decomposition $V = V_0 \oplus V_+ \oplus V_-$ satisfying (D.13).

By Theorem D.11, we know that $V$ has an orthogonal basis $B = \{v_1, \ldots, v_n\}$. Reordering this basis if necessary, we may assume that

\begin{align*}
B(v_i, v_i) &= 0, \quad i = 1, \ldots, r \\
B(v_i, v_i) &> 0, \quad i = r + 1, \ldots, p \\
B(v_i, v_i) &< 0, \quad i = p, \ldots, q.
\end{align*}

So $V_0 = \bigoplus_{i=1}^{r}$, $V_+ = \bigoplus_{i=r+1}^{p}$, $V_- = \bigoplus_{i=p+1}^{q}$ work.

Let $V = V_0 \oplus V_+ \oplus V_-$ satisfy (D.13) and $W \subset V$ be a subspace satisfying $B(v,v) > 0$ for all $0 \neq v \in W$. Then we establish the following:

**Claim.** $W + V_0 + V_+ = W \oplus V_0 \oplus V_-.$

Let $w \in W$, $v_0 \in V_-$, $v_- \in V_-$ satisfy $w + v_0 + v_- = 0$. Then

\begin{align*}
0 &= B(w, w + v_0 + v_-) = B(w,w) + B(w,v_0) + B(w,v_-) \\
&= B(w,w) + B(w,v_-)
\end{align*}

As $B$ is symmetric, we have $B(v,v) = B(v,v)$. But $B(v,v) = 0$ if and only if $v = 0$ and $B(v,v) = 0$ if and only if $v = 0$. So $w = 0 = v_-$. Hence also $v_0 = 0$ and the claim is established.

The claim implies that $\dim W \leq \dim V_+$. In particular, if we have another orthogonal decomposition $V = V'_0 \oplus V'_+ \oplus V'_-$ satisfying (D.13), then $\dim V'_+ \leq \dim V_+$. Applying the same argument shows that $\dim V_+ \leq \dim V'_+$, so $\dim V_+ = \dim V'_+$. An analogous argument shows that $\dim V_- = \dim V'_-$, hence also that $\dim V_0 = \dim V'_0$.

**Remark D.14.** In the above the subspace $V_0 = \text{rad} V$, so is unique by Lemma D.6, but $V_+$ and $V_0$ are not, although their dimensions are.
**Definition D.15.** Let \((V, B)\) be a bilinear space over \(F\), a subfield of \(\mathbb{R}\) and \(<\) an ordering on \(F\) arising from \(\mathbb{R}\). Let \(V = V_0 \oplus V_+ \oplus V_-\) be an orthogonal decomposition satisfying (D.13). Then \(\dim V_+ - \dim V_-\) is called the signature of \(B\). We denote the signature by \(\text{sgn}_P(B)\), where \(P = \{x \in F \mid x > 0\}\).

In the above, if \(V = \text{rad}(B) \perp W\), then \((W, B|W \times W)\) is regular by Lemma D.6, unique up to isometry, so satisfies \(\text{sgn}_P(V) = \text{sgn}_P(W)\).

Proposition D.12 implies

**Corollary D.16.** Let \((V, B)\) be a bilinear space over \(F\), a subfield of \(\mathbb{R}\) and \(<\) an ordering on \(F\) arising from \(\mathbb{R}\). Then the signature of \(B\), \(\text{sgn}_P(B)\) is independent of orthogonal decomposition satisfying (D.13).

If \(F = \mathbb{R}\), then we classify isometry classes of finite dimensional real symmetric bilinear spaces.

**Theorem D.17.** (Sylvestor’s Law of Inertia) Two regular finite dimensional symmetric bilinear spaces over \(\mathbb{R}\) are isometric if and only if they have the same signature.

**Proof.** Let \(V = \mathbb{R}v\) be a one dimensional real vector space with symmetric bilinear form \(B\) given by \(B(v, v) > 0\). Then \(B(v, v) = x^2\), some \(x \in \mathbb{R}^\times\). The map \(V \to \mathbb{R}\) given by \(v \mapsto x\) determines an isometry \((V, B) \to (\mathbb{R}, B')\) with \(B'(1, 1) = 1\). Similarly, if \(B(v, v) = -x^2\), we get an isometry \((V, B) \to (\mathbb{R}, B')\) with \(B'(1, 1) = -1\). \(\square\)

**Exercise D.18.** Prove Lemma D.7.
APPENDIX E


Primitive Roots

Let $G$ be a finite abelian group. In this section, we compute the automorphism group of $G$. This is useful in the study of cyclotomic extensions of a field. We know that $G \cong \mathbb{Z}/n\mathbb{Z}$ for some positive integer $n$ and $\text{Aut}(\mathbb{Z}/n\mathbb{Z}) \cong (\mathbb{Z}/n\mathbb{Z})^\times$ (as generators must go to generators), so it is abelian of order $\varphi(n)$. By the Chinese Remainder Theorem, we are reduced to the case of computing $(\mathbb{Z}/p^n\mathbb{Z})^\times$ where $p$ is a (positive) prime. We also know that $(\mathbb{Z}/p\mathbb{Z})^\times$ is cyclic, so can assume that $n > 1$.

We shall show that $(\mathbb{Z}/p^n\mathbb{Z})^\times$ is also cyclic unless $p = 2$. A generator for this automorphism group is called a primitive root. The proofs rely on the binomial theorem and its special case, the Children’s Binomial Theorem: that $(a + b)^p \equiv a^p + b^p \mod p$ for any integers $a$ and $b$. In particular, it implies

**Lemma E.1.** Let $p$ be a (positive) prime and $n > 1$. If $a$, $b$, and $k$ are integers, then $b^p \equiv b^p + kp^n \mod p^{n+1}$. In particular, if $a \equiv b \mod p^n$, then $a^p \equiv b^p \mod p^{n+1}$.

**Lemma E.2.** Let $p$ be a (positive) odd prime and $n \geq 2$. Then for every integer $a$, we have $(1 + ap)^{p^{n-2}} \equiv 1 + ap^{n-1} \mod p^n$.

**Proof.** We induct on $n$. The case $n = 2$ is trivial, so we may assume that $n > 2$. In particular, $2(n-1) > n+1$, so $p(n-1) > n+1$. By induction, we may assume the result for $n-2$ and show the result for $n-1$. By the lemma, induction, and the Binomial Theorem, we have

$$
(1 + ap)^{p^{n-1}} \equiv ((1 + ap)^{p^{n-2}})^p \equiv (1 + ap^{n-1})^p \\
\equiv 1 + \binom{p}{1} ap^{n-1} + \binom{p}{2} a^2 p^{2(n-1)} + \cdots + (ap^{n-1})^p \\
\equiv 1 + \binom{p}{1} ap^{n-1} \equiv 1 + ap^n \mod p^{n+1}.
$$

**Lemma E.3.** Let $p$ be a (positive) odd prime and $a$ an integer not divisible by $p$. Then for every integer $n \geq 2$, the congruence class of $1 + ap$ in $\mathbb{Z}/p^n\mathbb{Z}$ has order $p^{n-1}$.
PROOF. Let $\overline{\pi} : \mathbb{Z} \to \mathbb{Z}/p^n\mathbb{Z}$ the canonical epimorphism. By the previous lemma, we know that $(1 + ap)^{p^{n-1}} \equiv 1 + ap^n \mod p^{n+1}$, hence $(1 + ap)^{p^{n-1}} \equiv 1 \mod p^n$ and the order of $(1 + ap)$ divides $p^{n-1}$. The previous lemma also implies that $(1 + ap)^{p^{n-1}} \equiv 1 + ap^{n-1} \not\equiv 1 \mod p^n$, so the order of $(1 + ap)$ must be $p^{n-1}$.

Proposition E.4. Let $p$ be a (positive) odd prime. Then $(\mathbb{Z}/p^n\mathbb{Z})^\times$ is cyclic for all $n$.

PROOF. Let $a$ be an integer relatively prime to $p$. We know that $(\mathbb{Z}/p\mathbb{Z})^\times$ is cyclic of order $p - 1$, so generated by the residue class of $a$.

Claim 1. We may assume that $a^{p-1} \not\equiv 1 \mod p^2$:

Suppose that $a^{p-1} \equiv 1 \mod p^2$. Let $b = a + p$. Then

$$b^{p-1} = (a + p)^{p-1} \equiv a^{p-1} + (p - 1)a^{p-2}p \equiv 1 + (p - 1)a^{p-2}p \not\equiv 1 \mod p^2$$

by the Binomial Theorem, so replacing $a$ by $b$ works.

Claim 2. Let $\overline{\pi} : \mathbb{Z} \to \mathbb{Z}/p^n\mathbb{Z}$ the canonical epimorphism. If $a$ is as in the first claim, then $\langle \overline{a} \rangle = (\mathbb{Z}/p^n\mathbb{Z})^\times$:

We know that $\varphi(p^n) = p^{n-1}(p - 1)$ is the order of $(\mathbb{Z}/p^n\mathbb{Z})^\times$, so it suffices to prove that if $\overline{a}^N = 1$, then $p^{n-1} | N$ and $p - 1 | N$ in $\mathbb{Z}$, as $p$ and $p - 1$ are relatively prime. By the first claim, we can write $a = 1 + xp$ for some integer $x$ not divisible by $p$. By the last lemma, $\overline{1 + xp}$ has order $p^{n-1}$ in $(\mathbb{Z}/p^n\mathbb{Z})^\times$. As $\overline{a}^N = 1$, we must have $p^{n-1} | N$. Write $N = p^{k-1}M$ for some integer $M$. By Fermat’s Little Theorem, $1 \equiv (a^N) = (a^M)^{p^{k-1}} \equiv a^M \mod p$, so $p - 1 \mid M$ as $a$ is relatively prime to $p$. Hence $\overline{a}^M$ has order $p - 1$. This establishes Claim 2 and the theorem. □

We know that $(\mathbb{Z}/2\mathbb{Z})^\times \cong 1$ and $(\mathbb{Z}/2^2\mathbb{Z})^\times \cong \mathbb{Z}/2\mathbb{Z}$ are cyclic, however, this is not the case for $(\mathbb{Z}/2^n\mathbb{Z})^\times$ when $n > 2$.

Proposition E.5. Let $n \geq 3$ be an integer. Then $(\mathbb{Z}/2^n\mathbb{Z})^\times \cong \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2^{n-2}\mathbb{Z}$.

PROOF. Let $n \geq 3$. We begin with the following:

Claim. $5^{2^{n-3}} \equiv 1 + 2^{n-1} \mod 2^n$:

This is true for $n = 3$, so we proceed by induction. As $n \geq 3$, we have $2n - 2 \geq n + 1$. Assuming the result for $n - 3$ we show it holds for $n - 2$. By Lemma E.1

$$5^{2^{n-2}} = (5^{2^{n-3}})^2 = (1 + 2^{n-1})^2 = 1 + 2^n + 2^{2n-2} \equiv 1 + 2^n \mod 2^{n+1},$$
proving the claim. The claim implies that $5^{2n-3} \equiv 1 + 2^{n-1} \not\equiv 1 \mod 2^n$ and $5^{2n-2} \equiv 1 + 2^n \mod 2^{n+1}$. It follows that $5^{2n-2} \equiv 1 \mod 2^n$, and the order of $\overline{5}$ in $(\mathbb{Z}/2^n\mathbb{Z})^\times$ is $2^{n-2}$, where $- : \mathbb{Z} \to \mathbb{Z}/2^n\mathbb{Z}$ is the canonical epimorphism. To finish, it suffices to show that the subgroup $\langle -1, \overline{5} \rangle$ in $\mathbb{Z}/2^n\mathbb{Z}$ has order $2^{n-1}$, as this would imply that $(\mathbb{Z}/2^n\mathbb{Z})^\times \cong \langle -1 \rangle \times \langle 5 \rangle$. So suppose that $(-1)^a \overline{5}^b \equiv (-1)^{a'} \overline{5}^{b'} \mod 2^n$ for some positive integers $a, a', b, b'$. As $n \geq 2$ and $5 \equiv 1 \mod 4$, we have $(-1)^a \equiv (-1)^{a'} \mod 4$, hence $a \equiv a' \mod 2$. It follows that $5^b \equiv 5^{b'} \mod 2^n$. We may assume that $b \geq b'$, hence $5^{b-b'} \equiv 1 \mod 2^n$. Since $\overline{5}$ has order $2^{n-2}$, we have $2^{n-2} \mid b - b'$, i.e., $b \equiv b' \mod 2^n$. The result follows.

The two propositions and the Chinese Remainder Theorem yield the desired result.

**Theorem E.6.** Let $n$ be a positive integer. Then $(\mathbb{Z}/n\mathbb{Z})^\times$ is cyclic if and only if $n = 2, 4, p^r$, or $2p^r$ where $p$ is an odd prime.
APPENDIX F

Pell’s Equation

In this section, we show that the well-known diophantine equation

(*) \[ x^2 - dy^2 = 1 \]

is solvable, i.e., there exist integers \( x, y \) satisfying (*). In elementary number theory, it is shown that the solutions may be found using continued fractions, a method based on repeated use of the division algorithm, although it may not so easy to calculate. We need two lemmas.

**Lemma F.1.** Let \( \alpha \) be an irrational number. Then there exist infinitely many relatively prime integers \( x, y \) satisfying

\[ \left| \frac{x}{y} - \alpha \right| < \frac{1}{y^2}. \]

**Proof.** Let \( n \) be a positive integer. Partition the half-open interval \([0, 1)\) into \( n \) half-open subintervals,

\[ [0, 1) = \bigcup_{j=0}^{n-1} \left[ \frac{j}{n}, \frac{j+1}{n} \right). \]

Recall if \( \beta \) is a real number then \([\beta]\) is defined to be the largest integer \( \leq \beta \), and \( \beta - [\beta] \) is then called the fractional part of \( \beta \). Consider the \( n + 1 \) fractional parts of the elements \( 0, \alpha, 2\alpha, \ldots, n\alpha \). By the Dirichlet Pigeon Hole Principle, two of these fractional parts must lie in the same subinterval in \((\dagger)\). Hence there exist integers \( j \) and \( k \) satisfying \( 0 \leq j < k \leq n \) and

\[ |(k\alpha - [k\alpha]) - (j\alpha - [j\alpha])| < \frac{1}{n}. \]

Set \( y = k - j \) and \( x = [j\alpha] - [k\alpha] \). We have \( 0 < y < n \) and

\[ 0 < \left| \frac{x}{y} - \alpha \right| < \frac{1}{yn} < \frac{1}{y^2}. \]

Therefore, \( x/y \) is one solution. We may assume that \( x \) and \( y \) are relatively prime, since dividing by the gcd of \( x \) and \( y \) would also lead to a solution. Now choose an integer \( n_1 > 1/|((x/y) - \alpha)| \). Repeating
the construction would yield a new solution \( x_1, y_1 \) in relatively prime integers. If follows that we can construct infinitely many such solutions.

\[ \square \]

**Lemma F.2.** Let \( d \) be a positive square-free integer. Then there exists a constant \( c \) satisfying

\[ |x^2 - dy^2| < c \]

has infinitely many solutions in integers \( x, y \).

**Proof.** We have \( x^2 - dy^2 = (x - y\sqrt{d})(x + y\sqrt{d}) \) in \( \mathbb{Z}[d] \). By the lemma, there exist infinitely many relatively prime integers \( x, y \) with \( y > 0 \) and satisfying \( |x - y\sqrt{d}| < 1/y \). Therefore,

\[ |x + y\sqrt{d}| < |x - y\sqrt{d}| + 2y\sqrt{d} < \frac{1}{y} + 2y\sqrt{d} \]

and

\[ |x^2 - dy^2| < \frac{1}{y^2} + 2\sqrt{d} < 1 + 2\sqrt{d}, \]

so \( c = 1 + 2\sqrt{d} \) works. \( \square \)

**Proposition F.3.** Let \( d \) be a square-free positive integer. Then Pell’s equation \( x^2 - dy^2 = 1 \) has infinitely many solutions in integers. Further, there exists an integral solution \((x_1, y_1)\) such that every solution is of the form \((x_n, y_n)\) with \( x_n + y_n\sqrt{d} = (x_1 + y_1\sqrt{d})^n \) for some integer \( n \). In particular, the solution set in integers to Pell’s equation forms an infinite cyclic group.

**Proof.** We begin with some notation. If \( \gamma = x + y\sqrt{d} \), \( x, y \) in \( \mathbb{Q} \), let \( \overline{\gamma} = x - y\sqrt{d} \) and \( N(\gamma) = \gamma\overline{\gamma} = x^2 - dy^2 \). By the last lemma, we know that there exists a nonzero integer \( m \) such that \( x^2 - dy^2 = m \) has infinitely many solutions \( x, y \) in positive integers. We may also assume that we have reduced this set of solutions to those having different \( x \)’s. In particular, among these infinitely many solutions in positive integers, there exist two such pairs \( x_1, y_1 \) and \( x_2, y_2 \) to this equation with \( x_1 \neq x_2, x_1 \equiv x_2 \mod |m| \), and \( y_1 \equiv y_2 \mod |m| \). Let \( \alpha = x_1 + y_1\sqrt{d} \) and \( \beta = x_2 + y_2\sqrt{d} \). By choice, \( N(\alpha) = m = N(\beta) \). Write \( \alpha\beta = r + s\sqrt{d} \) with \( r \) and \( s \) integers satisfying \( r \equiv 0 \equiv s \mod |m| \). It follows that

\[ \alpha\beta = m(a + b\sqrt{d}) \]

for some integers \( a \) and \( b \). Taking \( N \) of this equation yields \( a^2 - db^2 = 1 \). Suppose that \( b = 0 \), then \( a = \pm 1 \) and \( am = \alpha\beta\overline{\beta} = \pm m\beta \), i.e., \( a = \pm \beta \). This implies \( x_1 = x_2 \), a contradiction. It follows that Pell’s equation has a solution \( a, b \) to \( a^2 - db^2 = 1 \) with \( ab \neq 0 \).
Among all the solutions \(a, b\) to Pell’s equation choose one in positive integers \(a, b\) with the real number \(\alpha = a + b\sqrt{d} > 0\) minimal. Suppose that \(x, y\) is another solution in positive integers. Set \(\beta = x + y\sqrt{d}\). We show that \(\beta = \alpha^n\) for some positive integer \(n\). Suppose not. Then there exists a positive integer \(n\) such that \(\alpha^n < \beta < \alpha^{n+1}\). Since \(\alpha\overline{\alpha} = 1\), i.e., \(\alpha^{-1} = \overline{\alpha}\), we have

\[
1 < \beta\alpha^{-n} < \alpha.
\]

We can write \(1 < \beta\alpha^{-n} = \beta\overline{\alpha}^n = r + s\sqrt{d}\) with \(r\) and \(s\) integers. As \(N(\beta\overline{\alpha}^n) = N(\beta)(\overline{\alpha}) = 1\), we have \(r, s\) is also a solution of Pell’s equation. Since \(r + s\sqrt{d} > 0\), we have \(0 < r - s\sqrt{d} = (r + s\sqrt{d})^{-1} < 1\). In particular, \(r\) is positive and \(s\sqrt{d} > r - 1 \geq 0\). Therefore, \(s\) is also positive, contradicting the minimality of \(a, b\). If \(x, y\) is a solution in integers to Pell’s equation with \(x > 0\) and \(y < 0\), then \(\beta = x + y\sqrt{d}\) satisfies \(N(\beta) = 1\), so \(\beta = \alpha^n\) for some \(n\), hence \(\beta = \alpha^{-n}\). The cases \(x < 0, y > 0\) and \(x < 0, y < 0\) lead to solutions \(-\alpha^n\) for some \(n\) in \(\mathbb{Z}\). □
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Books on Field Theory


**Books on Number Theory**


**Book on Commutative Algebra**


**Books on Non-Commutative Algebra**


Book on with Material on Formally Real Fields


### Notation

\[
\begin{pmatrix}
1 & 2 & \cdots & n \\
\sigma(1) & \sigma(2) & \cdots & \sigma(n)
\end{pmatrix},
\]
representation of a permutation, 67

\((A_*, d^*),\) cochain complex, 251

\((A_*, d_*),\) chain complex, 251

\((L/F, f),\) crossed product algebra, 722

\((R, m),\) local ring with maximal ideal \(m,\) 608

\((\mathbb{Z}/m\mathbb{Z})^\times,\) units in \(\mathbb{Z}/m\mathbb{Z},\) 37

\((\mathfrak{M}; \mathfrak{M}),\) 622

\((a),\) principal ideal generated by \(a,\) 163

\((a_1 \cdots a_r), r\)-cycle, 143

\((a_1, \ldots, a_n),\) ideal generated by \(a_1, \ldots, a_n,\) 163

\((g_i)_i,\) elements in \(\times_i G_i,\) 52

\(*: G \times S \to S,\) G-action, 110

\(1_R,\) identity (unity) of \(R,\) 158

\(1_S,\) the identity map on \(S,\) 48

\(A < B,\) same as \(A \subset B, A \neq B,\) 62

\(A \setminus B,\) element in \(A\) not \(B,\) 3

\(A \xrightarrow{f} B,\) map from \(A\) to \(B,\) 7

\(A^{op},\) opposite algebra, 706

\(A^K,\) integral closure of \(A\) in \(K,\) 526

\(A_n,\) alternating group group on \(n\) letters, 67

\(A_n,\) alternating group on \(n\)-letters, 145

\(C(P_1, \ldots, P_n),\) constructible points from \(P_1, \ldots, P_n,\) 363

\(C(a),\) conjugacy class of \(a,\) 116

\(C(z_1, \ldots, z_n),\) constructible points from \(z_1, \ldots, z_n,\) 364

\(D(R),\) basic open set, 601

\(D_n,\) dihedral group of order \(2n,\) 50

\(D_n(R),\) diagonal group, 51

\(F(X),\) smallest field containing \(F\) and \(S,\) 337

\(F^S,\) fixed field of \(F\) under the set \(S\) of automorphisms, 381

\(F^p,\) \(\{x^p \mid x \in F\},\) 376

\(F_G(S),\) G-fixed point set, 112

\(G \cong G',\) isomorphic, 56

\(G(K/F),\) Galois group of \(K/F,\) 383

\(G/H,\) factor group of \(G\) by \(H,\) 72

\(G/H,\) set of cosets, 60

\(G^{(n)}, (G^{(n-1)})',\) 70, 95

\(G_{ab},\) abelianization of \(G,\) 80

\(G_\mathfrak{A}(R),\) 664

\(G_s,\) isotropy subgroup (stabilizer) of \(s,\) 111

\(H \rtimes \varphi G,\) semidirect product of \(H\) and \(G\) via \(\varphi,\) 69

\(H \triangleleft G,\) normal subgroup, 65

\(H \triangleleft\triangleleft G,\) \(H\) a characteristic subgroup of \(G,\) 68

\(H_n(A),\) homology of \((A_*, d_*)\), 251

\(K/F,\) field extension, 336

\(L_A(M);\) ring of left multiplications by \(A\) on \(M,\) 749

\(M_n(R),\) matrices coefficients in \(R,\) 36

\(N_G(H),\) normalizer of \(H\) in \(G,\) 53

\(Q,\) the rational numbers, 3

\(R,\) ring, 158

\(R,\) the set of real numbers, 7

\(R/\mathfrak{M},\) quotient (factor) ring of \(R\) by \(\mathfrak{M},\) 169

\(R[[t]],\) formal power series over \(R,\) 159

\(R^x,\) units in \(R,\) 47

\(R^x,\) units in ring \(R,\) 38

\(R^x,\) units of \(R,\) 158

\(R_p,\) localization at \(p,\) 605
\( R_a \), localization at set of the powers of \( a \), 605
\( S \lor T \), disjoint union of \( S \) and \( T \), 12
\( S^{-1}R \), localization of \( R \) by \( S \), 188
\( S_n \), symmetric group on \( n \) letters, 48
\( V_T(T) \), variety of \( T \), 285, 599
\( Z(G) \), center of \( G \), 66, 116
\( Z(R) \), center of \( R \), 162
\( Z_F(\mathfrak{A}) \), zero set in \( F^n \) of \( \mathfrak{A} \), 231
\( Z_G(a) \), centralizer of \( a \), 116
\( [G,G] = G' \), commutator subgroup, 70, 95
\( [G:H] \), index of \( H \) in \( G \), 61
\( [K:F] \), degree of \( K \) over \( F \), 336
\( [T]_{B,C} \), matrix representation relative to bases \( B,C \), 67
\( [x,y] \), commutator of \( x \) and \( y \), 95
\( [x]_c \), equivalence class of \( x \), 31
\( \text{Ass}_R V(\mathfrak{A}) \), associated primes of \( \mathfrak{A} \)
\( \text{Ass}_R V(\mathfrak{A}) \), associated primes of \( \mathfrak{A} \)
\( \text{Br}(F) \), Brauer group of \( F \), 718
\( \mathbb{C} \), set of complex numbers, 3
\( \text{char}(R) \), characteristic of \( R \), 171
\( \Delta_K \), prime subfield of \( K \), 176
\( \text{End}(A) \), endomorphism ring of \( A \), 168
\( \text{End}_F(V) \), endomorphism ring of \( F \)-vector space \( V \), 52
\( \text{Aut}_F(V) \), automorphism group of \( F \)-vector space \( V \), 122
\( \text{Br}(F) \), Brauer group of \( F \), 718
\( \mathbb{C} \), set of complex numbers, 3
\( \text{char}(R) \), characteristic of \( R \), 171
\( \Delta_K \), prime subfield of \( K \), 176
\( \text{End}(A) \), endomorphism ring of \( A \), 168
\( \text{End}_F(V) \), endomorphism ring of \( F \)-vector space \( V \), 52
\( \text{Ext}_R(M,N) \), cohomology from Hom, 259
\( \text{GL}_n(R) \), general linear group over \( R \), 51
\( \implies \), implies, 6
\( \mathbb{M}_n(R) \), \( n \times n \) matrix ring over \( R \), 30
\( \text{Max}(R) \), set of maximal ideals, 286
\( \text{Min}(R) \), set of minimal primes in \( R \), 602
\( \mathbb{Q}[t] \), polynomials with rational coefficients, 4
\( \mathbb{R}^+ \), the set of positive reals, 47
\( \Sigma(S) \), permutation group on \( S \), 48
\( \text{Spec}(R) \), Spectrum of \( R \), 285
\( \text{Syl}_p(G) \), set of Sylow \( p \)-subgroups of \( G \), 120, 126
\( \text{Tor}^R(M,N) \), homology from \( \otimes_R \), 271
\( \mathbb{Z} \), the set of integers, 3
\( \mathbb{Z}/m\mathbb{Z} \), integers modulo \( m \), 34
\( \mathbb{Z}^+ \), non-negative integers, 7
\( \mathbb{Z}^+ \), positive integers, 7
\( \mathbb{Z}^+ \), the set of positive integers, 9
\( \approx \), usually an equivalence relation, 30
\( \bigcup I A_i \), the union of the sets \( A_i \), 32
\( \bigvee I A_i \), disjoint union of the sets \( i \) in \( I \), 32
\( \mathbb{C}^* \), is \( = \mathbb{C} \setminus \{0\} \), 116
\( \mathcal{F}(K/F) \), intermediate fields of \( K/F \), 403
\( \mathcal{G} \), set of subgroups of \( G \), 120
\( \mathcal{G}(K/F) \), subgroups of \( G(K/F) \), 403
\( \mathcal{O} \), system of representatives for a \( G \)-action, 110
\( \mathcal{O}^* \), is \( \mathcal{O} \setminus \{0\} \), 112
\( \mathcal{P}(S) \), power set of \( S \), 119
\( \cong \), isomorphism, 161
\( \text{deg } A \), index of the algebra \( A \), 711
\( \text{det} \), determinant, 57
\( \emptyset \), the empty set, 11
\( \equiv \mod m \), congruence modulo \( m \), 34
\( \mathfrak{A}, \mathfrak{B}, \mathfrak{C} \), ideals, 162
\( \mathfrak{A}\mathfrak{B} \), ideal generated by \( \mathfrak{A} \) and \( \mathfrak{B} \), 164
\( \mathfrak{A} | \mathfrak{B} \), ideal division, 535
\( \mathfrak{m} \), maximal ideal, 166
\( \varphi \), \( \mathfrak{P} \), prime ideals, 167
\( \text{im } \varphi \), image of \( \varphi \), 56, 160
\( \text{ind}_H^G(W) \), induced module, 773
\( \text{ker } \varphi \), kernel of \( \varphi \), 56, 160
\( \langle W \rangle \), group (etc.) generated by \( W \), 49
\( \langle a \rangle \), group (etc) generated by \( a \), 49
\( \lambda_a \), left translation by \( a \), 61
\( \lambda_x \), left multiplication by \( x \), 73
\( | A | \), cardinality of \( A \), 6
\( \mu_n \), \( n \)-th roots of unity \( (\text{in } \mathbb{C}) \), 49
\( \text{nil}(R) \), nilradical of \( R \), 185
\( \text{SL}_n(R) \), special linear group over \( R \), 51
\( \overline{A} \), set of equivalence class of \( A \), 31
\( \varphi \) : \( G \rightarrow G/H \), canonical group epimorphism, 60
\( \pi \), equivalence class of \( x \), 31
\( a \varphi \), 603
\( \pi(x) \), 3
\( \prod R_i \times R_i \), product of rings, 176
\( \text{rad}(R) \), Jacobson radical of \( R \), 613
\( \text{res}_{G}^{H}(V) \), restriction module, 772
\( \text{sgn} \), signum (sign) map, 146
\( \sim \), usually an equivalence relation, 30
\( \sim_{G} \), \( G \)=equivalence, 110
\( \sqrt{A} \), radical of \( A \), 187, 601
\( \theta_x \), conjugation by \( x \), 65
\( \varphi : R[t] \rightarrow S[t] \), ring extension of
\( \varphi : R \rightarrow S \), 350
\( \text{tr deg}_F A \), \( \text{tr deg}_F qf(A) \), 646
\( \text{tr deg}_F K \), transcendence degree of \( K/F \), 348
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\( \text{zd}(R) \), zero divisors of \( R \), 183
\( a \approx b \), \( a \) is an associate of \( b \), 165
\( a \mid b \), \( a \) divides \( b \), 3, 165
\( a \not| b \), \( a \) does not divide \( b \), 3
\( a^{-1} \), inverse of \( a \), 46
\( e = e_G \), unity in the group \( G \), 46
\( f : A \rightarrow B \), map from \( A \) to \( B \), 6
\( f^{-1}(D) \), preimage of \( D \), 40, 75
\( f^{-1}(b) \), fiber of \( f \) at \( b \), 40
\( l(M) \), length of a module \( M \) having a composition series., 636
\( m_F(\alpha) \), minimal polynomial of \( \alpha \) over \( F \), 340
\( qf(R) \), quotient field of domain \( R \), 173
\( t \), represents a variable, 4
\( v_p(\mathfrak{A}) \), power of \( p \) dividing \( \mathfrak{A} \), 535
\( xHx^{-1} \), conjugate of \( H \), 65
\( \text{Sp}_n(R) \), symplectic group over \( R \), 51
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