Abstract Algebra

Theory and Applications

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Sage Exercises for Abstract Algebra

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Preface

This text is intended for a one or two-semester undergraduate course in abstract algebra. Traditionally, these courses have covered the theoretical aspects of groups, rings, and fields. However, with the development of computing in the last several decades, applications that involve abstract algebra and discrete mathematics have become increasingly important, and many science, engineering, and computer science students are now electing to minor in mathematics. Though theory still occupies a central role in the subject of abstract algebra and no student should go through such a course without a good notion of what a proof is, the importance of applications such as coding theory and cryptography has grown significantly.

Until recently most abstract algebra texts included few if any applications. However, one of the major problems in teaching an abstract algebra course is that for many students it is their first encounter with an environment that requires them to do rigorous proofs. Such students often find it hard to see the use of learning to prove theorems and propositions; applied examples help the instructor provide motivation.

This text contains more material than can possibly be covered in a single semester. Certainly there is adequate material for a two-semester course, and perhaps more; however, for a one-semester course it would be quite easy to omit selected chapters and still have a useful text. The order of presentation of topics is standard: groups, then rings, and finally fields. Emphasis can be placed either on theory or on applications. A typical one-semester course might cover groups and rings while briefly touching on field theory, using Chapters 1 through 6, 9, 10, 11, 13 (the first part), 16, 17, 18 (the first part), 20, and 21. Parts of these chapters could be deleted and applications substituted according to the interests of the students and the instructor. A two-semester course emphasizing theory might cover Chapters 1 through 6, 9, 10, 11, 13 through 18, 20, 21, 22 (the first part), and 23. On the other hand, if applications are to be emphasized, the course might cover Chapters 1 through 14, and 16 through 22. In an applied course, some of the more theoretical results could be assumed or omitted. A chapter dependency chart appears below. (A broken line indicates a partial dependency.)

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            ┌───────────┤ Chapters 1–6 ├───────────┐
            │           └──────┬───────┘           │
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┄┄┄┄┄┤ Chapter 8 │      │ Chapter 9 │       │ Chapter 7 │
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┄─┤ Chapter 13 │     │ Chapter 16 │   │ Chapter 12 │ │ Chapter 14 │
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│         │          ┌─────────┴──┐                  ┌───────┴────┐
│         │          │ Chapter 17 ├───┐              │ Chapter 15 │
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│  │ Chapter 18 │   │ Chapter 20 │  │ Chapter 19 │           ┆
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│                  ┌───────┴────┐                            ┆
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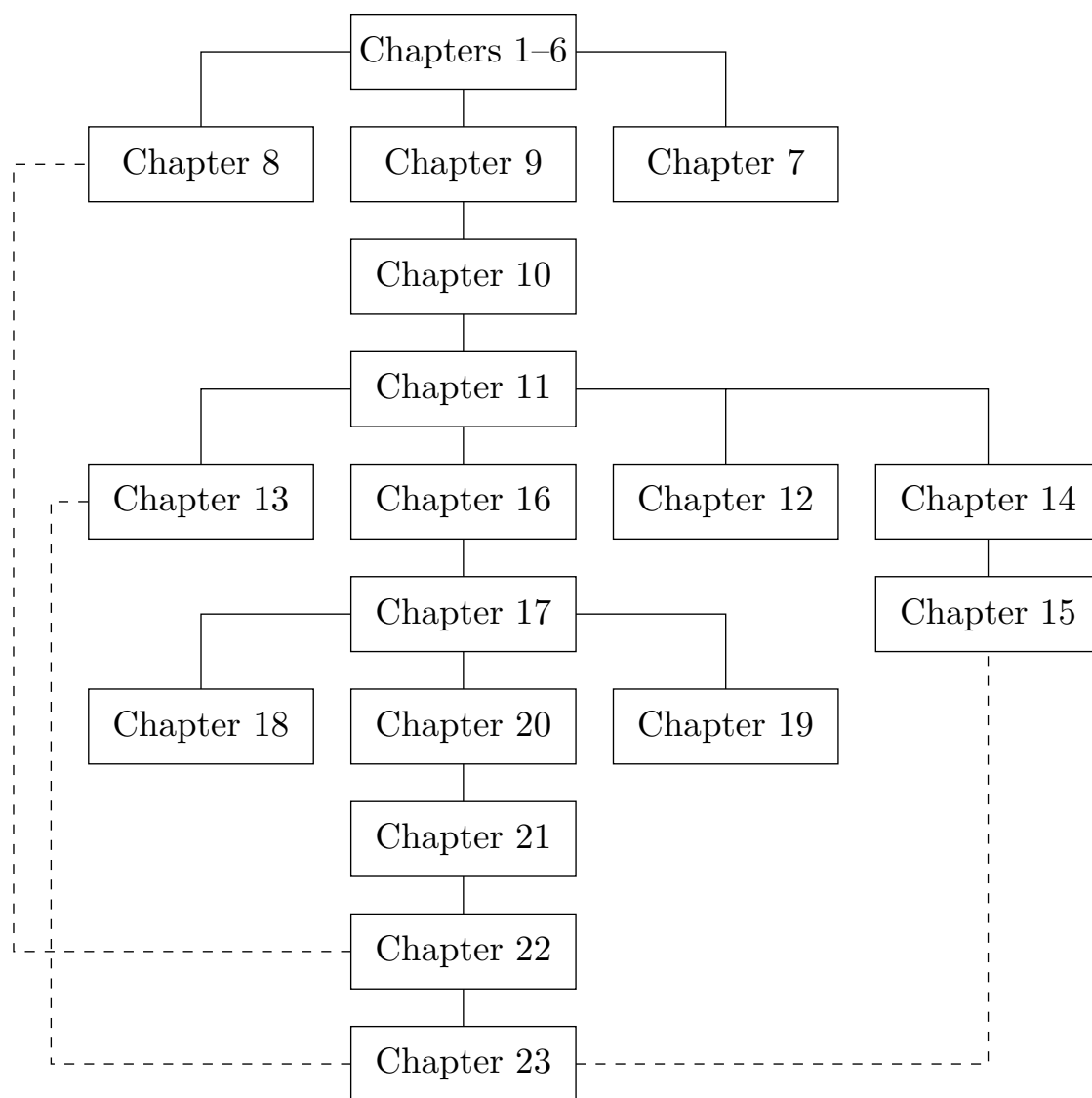
Though there are no specific prerequisites for a course in abstract algebra, students who have had other higher-level courses in mathematics will generally be more prepared than those who have not, because they will possess a bit more mathematical sophistication. Occasionally, we shall assume some basic linear algebra; that is, we shall take for granted an elementary knowledge of matrices and determinants. This should present no great problem, since most students taking a course in abstract algebra have been introduced to matrices and determinants elsewhere in their career, if they have not already taken a sophomore or junior-level course in linear algebra.

Exercise sections are the heart of any mathematics text. An exercise set appears at the end of each chapter. The nature of the exercises ranges over several categories; computational, conceptual, and theoretical problems are included. A section presenting hints and solutions to many of the exercises appears at the end of the text. Often in the solutions a proof is only sketched, and it is up to the student to provide the details. The exercises range in difficulty from very easy to very challenging. Many of the more substantial problems require careful thought, so the student should not be discouraged if the solution is not forthcoming after a few minutes of work.

There are additional exercises or computer projects at the ends of many of the chapters. The computer projects usually require a knowledge of programming. All of these exercises

and projects are more substantial in nature and allow the exploration of new results and theory.

Sage (sagemath.org) is a free, open source, software system for advanced mathematics, which is ideal for assisting with a study of abstract algebra. Sage can be used either on your own computer, a local server, or on CoCalc (cocalc.com). Robert Beezer has written a comprehensive introduction to Sage and a selection of relevant exercises that appear at the end of each chapter, including live Sage cells in the web version of the book. All of the Sage code has been subject to automated tests of accuracy, using the most recent version available at this time: SageMath Version 8.2 (released 2018-05-05).

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Contents

Preliminaries

A certain amount of mathematical maturity is necessary to find and study applications of abstract algebra. A basic knowledge of set theory, mathematical induction, equivalence relations, and matrices is a must. Even more important is the ability to read and understand mathematical proofs. In this chapter we will outline the background needed for a course in abstract algebra.

1.1   A Short Note on Proofs

Abstract mathematics is different from other sciences. In laboratory sciences such as chemistry and physics, scientists perform experiments to discover new principles and verify theories. Although mathematics is often motivated by physical experimentation or by computer simulations, it is made rigorous through the use of logical arguments. In studying abstract mathematics, we take what is called an axiomatic approach; that is, we take a collection of objects $\mathcal{S}$ and assume some rules about their structure. These rules are called ***axioms***. Using the axioms for $\mathcal{S}$, we wish to derive other information about $\mathcal{S}$ by using logical arguments. We require that our axioms be consistent; that is, they should not contradict one another. We also demand that there not be too many axioms. If a system of axioms is too restrictive, there will be few examples of the mathematical structure.

A ***statement*** in logic or mathematics is an assertion that is either true or false. Consider the following examples:

- $3 + 56 - 13 + 8/2$.

- All cats are black.

- $2 + 3 = 5$.

- $2x = 6$ exactly when $x = 4$.

- If $ax^2 + bx + c = 0$ and $a \neq 0$, then

$$x = \frac{-b \pm \sqrt{b^2 - 4ac}}{2a}.$$

- $x^3 - 4x^2 + 5x - 6$.

All but the first and last examples are statements, and must be either true or false.

A ***mathematical proof*** is nothing more than a convincing argument about the accuracy of a statement. Such an argument should contain enough detail to convince the audience; for instance, we can see that the statement "$2x = 6$ exactly when $x = 4$" is false by evaluating $2 \cdot 4$ and noting that $6 \neq 8$, an argument that would satisfy anyone. Of course, audiences

may vary widely: proofs can be addressed to another student, to a professor, or to the reader of a text. If more detail than needed is presented in the proof, then the explanation will be either long-winded or poorly written. If too much detail is omitted, then the proof may not be convincing. Again it is important to keep the audience in mind. High school students require much more detail than do graduate students. A good rule of thumb for an argument in an introductory abstract algebra course is that it should be written to convince one's peers, whether those peers be other students or other readers of the text.

Let us examine different types of statements. A statement could be as simple as "$10/5 = 2$;" however, mathematicians are usually interested in more complex statements such as "If p, then q," where p and q are both statements. If certain statements are known or assumed to be true, we wish to know what we can say about other statements. Here p is called the **hypothesis** and q is known as the **conclusion**. Consider the following statement: If $ax^2 + bx + c = 0$ and $a \neq 0$, then

$$x = \frac{-b \pm \sqrt{b^2 - 4ac}}{2a}.$$

The hypothesis is $ax^2 + bx + c = 0$ and $a \neq 0$; the conclusion is

$$x = \frac{-b \pm \sqrt{b^2 - 4ac}}{2a}.$$

Notice that the statement says nothing about whether or not the hypothesis is true. However, if this entire statement is true and we can show that $ax^2 + bx + c = 0$ with $a \neq 0$ is true, then the conclusion *must* be true. A proof of this statement might simply be a series of equations:

$$ax^2 + bx + c = 0$$
$$x^2 + \frac{b}{a}x = -\frac{c}{a}$$
$$x^2 + \frac{b}{a}x + \left(\frac{b}{2a}\right)^2 = \left(\frac{b}{2a}\right)^2 - \frac{c}{a}$$
$$\left(x + \frac{b}{2a}\right)^2 = \frac{b^2 - 4ac}{4a^2}$$
$$x + \frac{b}{2a} = \frac{\pm\sqrt{b^2 - 4ac}}{2a}$$
$$x = \frac{-b \pm \sqrt{b^2 - 4ac}}{2a}.$$

If we can prove a statement true, then that statement is called a **proposition**. A proposition of major importance is called a **theorem**. Sometimes instead of proving a theorem or proposition all at once, we break the proof down into modules; that is, we prove several supporting propositions, which are called **lemmas**, and use the results of these propositions to prove the main result. If we can prove a proposition or a theorem, we will often, with very little effort, be able to derive other related propositions called **corollaries**.

Some Cautions and Suggestions

There are several different strategies for proving propositions. In addition to using different methods of proof, students often make some common mistakes when they are first learning how to prove theorems. To aid students who are studying abstract mathematics for the

first time, we list here some of the difficulties that they may encounter and some of the strategies of proof available to them. It is a good idea to keep referring back to this list as a reminder. (Other techniques of proof will become apparent throughout this chapter and the remainder of the text.)

- A theorem cannot be proved by example; however, the standard way to show that a statement is not a theorem is to provide a counterexample.

- Quantifiers are important. Words and phrases such as *only*, *for all*, *for every*, and *for some* possess different meanings.

- Never assume any hypothesis that is not explicitly stated in the theorem. *You cannot take things for granted.*

- Suppose you wish to show that an object *exists* and is *unique*. First show that there actually is such an object. To show that it is unique, assume that there are two such objects, say r and s, and then show that $r = s$.

- Sometimes it is easier to prove the contrapositive of a statement. Proving the statement "If p, then q" is exactly the same as proving the statement "If not q, then not p."

- Although it is usually better to find a direct proof of a theorem, this task can sometimes be difficult. It may be easier to assume that the theorem that you are trying to prove is false, and to hope that in the course of your argument you are forced to make some statement that cannot possibly be true.

Remember that one of the main objectives of higher mathematics is proving theorems. Theorems are tools that make new and productive applications of mathematics possible. We use examples to give insight into existing theorems and to foster intuitions as to what new theorems might be true. Applications, examples, and proofs are tightly interconnected— much more so than they may seem at first appearance.

1.2 Sets and Equivalence Relations

Set Theory

A **set** is a well-defined collection of objects; that is, it is defined in such a manner that we can determine for any given object x whether or not x belongs to the set. The objects that belong to a set are called its **elements** or **members**. We will denote sets by capital letters, such as A or X; if a is an element of the set A, we write $a \in A$.

A set is usually specified either by listing all of its elements inside a pair of braces or by stating the property that determines whether or not an object x belongs to the set. We might write

$$X = \{x_1, x_2, \ldots, x_n\}$$

for a set containing elements $x_1, x_2, \ldots, x_n$ or

$$X = \{x : x \text{ satisfies } \mathcal{P}\}$$

if each x in X satisfies a certain property $\mathcal{P}$. For example, if E is the set of even positive integers, we can describe E by writing either

$$E = \{2, 4, 6, \ldots\} \quad \text{or} \quad E = \{x : x \text{ is an even integer and } x > 0\}.$$

We write $2 \in E$ when we want to say that 2 is in the set E, and $-3 \notin E$ to say that -3 is not in the set E.

Some of the more important sets that we will consider are the following:

$$\mathbb{N} = \{n : n \text{ is a natural number}\} = \{1, 2, 3, \ldots\};$$
$$\mathbb{Z} = \{n : n \text{ is an integer}\} = \{\ldots, -1, 0, 1, 2, \ldots\};$$
$$\mathbb{Q} = \{r : r \text{ is a rational number}\} = \{p/q : p, q \in \mathbb{Z} \text{ where } q \neq 0\};$$
$$\mathbb{R} = \{x : x \text{ is a real number}\};$$
$$\mathbb{C} = \{z : z \text{ is a complex number}\}.$$

We can find various relations between sets as well as perform operations on sets. A set A is a **subset** of B, written $A \subset B$ or $B \supset A$, if every element of A is also an element of B. For example,

$$\{4, 5, 8\} \subset \{2, 3, 4, 5, 6, 7, 8, 9\}$$

and

$$\mathbb{N} \subset \mathbb{Z} \subset \mathbb{Q} \subset \mathbb{R} \subset \mathbb{C}.$$

Trivially, every set is a subset of itself. A set B is a **proper subset** of a set A if $B \subset A$ but $B \neq A$. If A is not a subset of B, we write $A \not\subset B$; for example, $\{4, 7, 9\} \not\subset \{2, 4, 5, 8, 9\}$. Two sets are **equal**, written $A = B$, if we can show that $A \subset B$ and $B \subset A$.

It is convenient to have a set with no elements in it. This set is called the **empty set** and is denoted by $\emptyset$. Note that the empty set is a subset of every set.

To construct new sets out of old sets, we can perform certain operations: the **union** $A \cup B$ of two sets A and B is defined as

$$A \cup B = \{x : x \in A \text{ or } x \in B\};$$

the **intersection** of A and B is defined by

$$A \cap B = \{x : x \in A \text{ and } x \in B\}.$$

If $A = \{1, 3, 5\}$ and $B = \{1, 2, 3, 9\}$, then

$$A \cup B = \{1, 2, 3, 5, 9\} \quad \text{and} \quad A \cap B = \{1, 3\}.$$

We can consider the union and the intersection of more than two sets. In this case we write

$$\bigcup_{i=1}^{n} A_i = A_1 \cup \ldots \cup A_n$$

and

$$\bigcap_{i=1}^{n} A_i = A_1 \cap \ldots \cap A_n$$

for the union and intersection, respectively, of the sets $A_1, \ldots, A_n$.

When two sets have no elements in common, they are said to be **disjoint**; for example, if E is the set of even integers and O is the set of odd integers, then E and O are disjoint. Two sets A and B are disjoint exactly when $A \cap B = \emptyset$.

Sometimes we will work within one fixed set U, called the **universal set**. For any set $A \subset U$, we define the **complement** of A, denoted by A', to be the set

$$A' = \{x : x \in U \text{ and } x \notin A\}.$$

We define the **difference** of two sets A and B to be

$$A \setminus B = A \cap B' = \{x : x \in A \text{ and } x \notin B\}.$$

Example 1.1. Let $\mathbb{R}$ be the universal set and suppose that

$$A = \{x \in \mathbb{R} : 0 < x \leq 3\} \quad \text{and} \quad B = \{x \in \mathbb{R} : 2 \leq x < 4\}.$$

Then

$$\begin{aligned}
A \cap B &= \{x \in \mathbb{R} : 2 \leq x \leq 3\} \\
A \cup B &= \{x \in \mathbb{R} : 0 < x < 4\} \\
A \setminus B &= \{x \in \mathbb{R} : 0 < x < 2\} \\
A' &= \{x \in \mathbb{R} : x \leq 0 \text{ or } x > 3\}.
\end{aligned}$$

Proposition 1.2. *Let A, B, and C be sets. Then*

1. *$A \cup A = A$, $A \cap A = A$, and $A \setminus A = \emptyset$;*

2. *$A \cup \emptyset = A$ and $A \cap \emptyset = \emptyset$;*

3. *$A \cup (B \cup C) = (A \cup B) \cup C$ and $A \cap (B \cap C) = (A \cap B) \cap C$;*

4. *$A \cup B = B \cup A$ and $A \cap B = B \cap A$;*

5. *$A \cup (B \cap C) = (A \cup B) \cap (A \cup C)$;*

6. *$A \cap (B \cup C) = (A \cap B) \cup (A \cap C)$.*

PROOF.   We will prove (1) and (3) and leave the remaining results to be proven in the exercises.

(1) Observe that

$$\begin{aligned}
A \cup A &= \{x : x \in A \text{ or } x \in A\} \\
&= \{x : x \in A\} \\
&= A
\end{aligned}$$

and

$$\begin{aligned}
A \cap A &= \{x : x \in A \text{ and } x \in A\} \\
&= \{x : x \in A\} \\
&= A.
\end{aligned}$$

Also, $A \setminus A = A \cap A' = \emptyset$.

(3) For sets A, B, and C,

$$\begin{aligned}
A \cup (B \cup C) &= A \cup \{x : x \in B \text{ or } x \in C\} \\
&= \{x : x \in A \text{ or } x \in B, \text{ or } x \in C\} \\
&= \{x : x \in A \text{ or } x \in B\} \cup C \\
&= (A \cup B) \cup C.
\end{aligned}$$

A similar argument proves that $A \cap (B \cap C) = (A \cap B) \cap C$. $\square$

Theorem 1.3 De Morgan's Laws. *Let A and B be sets. Then*

1. *$(A \cup B)' = A' \cap B'$;*

2. *$(A \cap B)' = A' \cup B'$.*

PROOF.   (1) If $A \cup B = \emptyset$, then the theorem follows immediately since both A and B are the empty set. Otherwise, we must show that $(A \cup B)' \subset A' \cap B'$ and $(A \cup B)' \supset A' \cap B'$. Let $x \in (A \cup B)'$. Then $x \notin A \cup B$. So x is neither in A nor in B, by the definition of the union of sets. By the definition of the complement, $x \in A'$ and $x \in B'$. Therefore, $x \in A' \cap B'$ and we have $(A \cup B)' \subset A' \cap B'$.

To show the reverse inclusion, suppose that $x \in A' \cap B'$. Then $x \in A'$ and $x \in B'$, and so $x \notin A$ and $x \notin B$. Thus $x \notin A \cup B$ and so $x \in (A \cup B)'$. Hence, $(A \cup B)' \supset A' \cap B'$ and so $(A \cup B)' = A' \cap B'$.

The proof of (2) is left as an exercise. $\square$

Example 1.4. Other relations between sets often hold true. For example,

$$(A \setminus B) \cap (B \setminus A) = \emptyset.$$

To see that this is true, observe that

$$\begin{aligned}
(A \setminus B) \cap (B \setminus A) &= (A \cap B') \cap (B \cap A') \\
&= A \cap A' \cap B \cap B' \\
&= \emptyset.
\end{aligned}$$

Cartesian Products and Mappings

Given sets A and B, we can define a new set $A \times B$, called the **Cartesian product** of A and B, as a set of ordered pairs. That is,

$$A \times B = \{(a,b) : a \in A \text{ and } b \in B\}.$$

Example 1.5. If $A = \{x, y\}$, $B = \{1, 2, 3\}$, and $C = \emptyset$, then $A \times B$ is the set

$$\{(x, 1), (x, 2), (x, 3), (y, 1), (y, 2), (y, 3)\}$$

and

$$A \times C = \emptyset.$$

We define the **Cartesian product of n sets** to be

$$A_1 \times \cdots \times A_n = \{(a_1, \ldots, a_n) : a_i \in A_i \text{ for } i = 1, \ldots, n\}.$$

If $A = A_1 = A_2 = \cdots = A_n$, we often write A^n for $A \times \cdots \times A$ (where A would be written n times). For example, the set $\mathbb{R}^3$ consists of all of 3-tuples of real numbers.

Subsets of $A \times B$ are called **relations**. We will define a **mapping** or **function** $f \subset A \times B$ from a set A to a set B to be the special type of relation where $(a, b) \in f$ if for every element $a \in A$ there exists a unique element $b \in B$. Another way of saying this is that for every element in A, f assigns a unique element in B. We usually write $f : A \to B$ or $A \xrightarrow{f} B$. Instead of writing down ordered pairs $(a, b) \in A \times B$, we write $f(a) = b$ or $f : a \mapsto b$. The set A is called the **domain** of f and

$$f(A) = \{f(a) : a \in A\} \subset B$$

is called the **range** or **image** of f. We can think of the elements in the function's domain as input values and the elements in the function's range as output values.
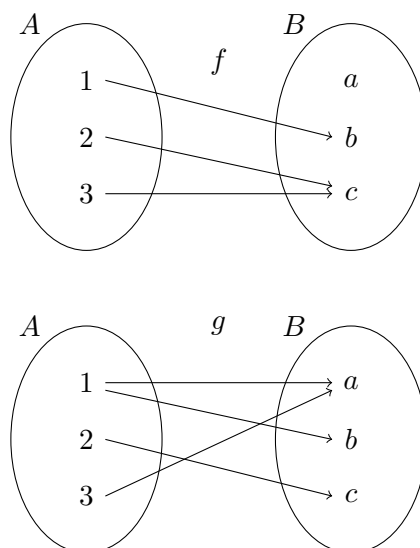
Figure 1.6: Mappings and relations

Example 1.7. Suppose $A = \{1, 2, 3\}$ and $B = \{a, b, c\}$. In Figure 1.6 we define relations f and g from A to B. The relation f is a mapping, but g is not because $1 \in A$ is not assigned to a unique element in B; that is, $g(1) = a$ and $g(1) = b$.

Given a function $f : A \to B$, it is often possible to write a list describing what the function does to each specific element in the domain. However, not all functions can be described in this manner. For example, the function $f : \mathbb{R} \to \mathbb{R}$ that sends each real number to its cube is a mapping that must be described by writing $f(x) = x^3$ or $f : x \mapsto x^3$.

Consider the relation $f : \mathbb{Q} \to \mathbb{Z}$ given by $f(p/q) = p$. We know that $1/2 = 2/4$, but is $f(1/2) = 1$ or 2? This relation cannot be a mapping because it is not well-defined. A relation is **well-defined** if each element in the domain is assigned to a *unique* element in the range.

If $f : A \to B$ is a map and the image of f is B, i.e., $f(A) = B$, then f is said to be **onto** or **surjective**. In other words, if there exists an $a \in A$ for each $b \in B$ such that $f(a) = b$, then f is onto. A map is **one-to-one** or **injective** if $a_1 \neq a_2$ implies $f(a_1) \neq f(a_2)$. Equivalently, a function is one-to-one if $f(a_1) = f(a_2)$ implies $a_1 = a_2$. A map that is both one-to-one and onto is called **bijective**.

Example 1.8. Let $f : \mathbb{Z} \to \mathbb{Q}$ be defined by $f(n) = n/1$. Then f is one-to-one but not onto. Define $g : \mathbb{Q} \to \mathbb{Z}$ by $g(p/q) = p$ where p/q is a rational number expressed in its lowest terms with a positive denominator. The function g is onto but not one-to-one.

Given two functions, we can construct a new function by using the range of the first function as the domain of the second function. Let $f : A \to B$ and $g : B \to C$ be mappings. Define a new map, the **composition** of f and g from A to C, by $(g \circ f)(x) = g(f(x))$.
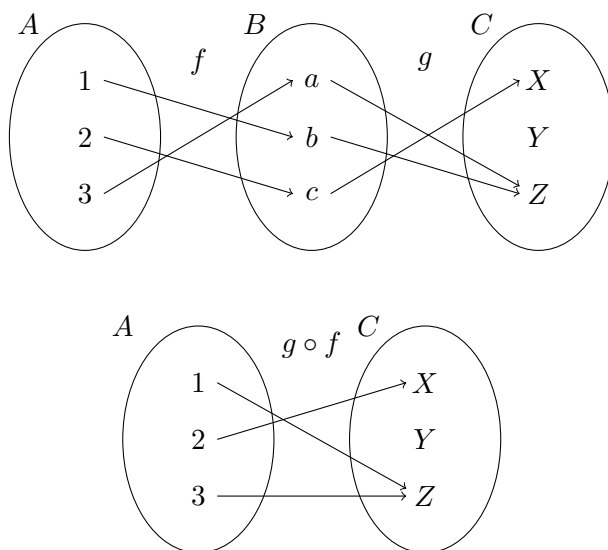
Figure 1.9: Composition of maps

Example 1.10. Consider the functions $f : A \to B$ and $g : B \to C$ that are defined in Figure 1.9 (top). The composition of these functions, $g \circ f : A \to C$, is defined in Figure 1.9 (bottom).

Example 1.11. Let $f(x) = x^2$ and $g(x) = 2x + 5$. Then

$$(f \circ g)(x) = f(g(x)) = (2x + 5)^2 = 4x^2 + 20x + 25$$

and

$$(g \circ f)(x) = g(f(x)) = 2x^2 + 5.$$

In general, order makes a difference; that is, in most cases $f \circ g \neq g \circ f$.

Example 1.12. Sometimes it is the case that $f \circ g = g \circ f$. Let $f(x) = x^3$ and $g(x) = \sqrt[3]{x}$. Then

$$(f \circ g)(x) = f(g(x)) = f(\sqrt[3]{x}) = (\sqrt[3]{x})^3 = x$$

and

$$(g \circ f)(x) = g(f(x)) = g(x^3) = \sqrt[3]{x^3} = x.$$

Example 1.13. Given a 2×2 matrix

$$A = \begin{pmatrix} a & b \\ c & d \end{pmatrix},$$

we can define a map $T_A : \mathbb{R}^2 \to \mathbb{R}^2$ by

$$T_A(x, y) = (ax + by, cx + dy)$$

for (x, y) in $\mathbb{R}^2$. This is actually matrix multiplication; that is,

$$\begin{pmatrix} a & b \\ c & d \end{pmatrix} \begin{pmatrix} x \\ y \end{pmatrix} = \begin{pmatrix} ax + by \\ cx + dy \end{pmatrix}.$$

Maps from $\mathbb{R}^n$ to $\mathbb{R}^m$ given by matrices are called ***linear maps*** or ***linear transformations***.

Example 1.14. Suppose that $S = \{1, 2, 3\}$. Define a map $\pi : S \to S$ by

$$\pi(1) = 2, \qquad \pi(2) = 1, \qquad \pi(3) = 3.$$

This is a bijective map. An alternative way to write π is

$$\begin{pmatrix} 1 & 2 & 3 \\ \pi(1) & \pi(2) & \pi(3) \end{pmatrix} = \begin{pmatrix} 1 & 2 & 3 \\ 2 & 1 & 3 \end{pmatrix}.$$

For any set S, a one-to-one and onto mapping $\pi : S \to S$ is called a **permutation** of S.

Theorem 1.15. *Let $f : A \to B$, $g : B \to C$, and $h : C \to D$. Then*

1. *The composition of mappings is associative; that is, $(h \circ g) \circ f = h \circ (g \circ f)$;*

2. *If f and g are both one-to-one, then the mapping $g \circ f$ is one-to-one;*

3. *If f and g are both onto, then the mapping $g \circ f$ is onto;*

4. *If f and g are bijective, then so is $g \circ f$.*

PROOF. We will prove (1) and (3). Part (2) is left as an exercise. Part (4) follows directly from (2) and (3).

(1) We must show that
$$h \circ (g \circ f) = (h \circ g) \circ f.$$

For $a \in A$ we have

$$\begin{aligned}
(h \circ (g \circ f))(a) &= h((g \circ f)(a)) \\
&= h(g(f(a))) \\
&= (h \circ g)(f(a)) \\
&= ((h \circ g) \circ f)(a).
\end{aligned}$$

(3) Assume that f and g are both onto functions. Given $c \in C$, we must show that there exists an $a \in A$ such that $(g \circ f)(a) = g(f(a)) = c$. However, since g is onto, there is an element $b \in B$ such that $g(b) = c$. Similarly, there is an $a \in A$ such that $f(a) = b$. Accordingly,
$$(g \circ f)(a) = g(f(a)) = g(b) = c. \qquad \square$$

If S is any set, we will use id_S or id to denote the **identity mapping** from S to itself. Define this map by $id(s) = s$ for all $s \in S$. A map $g : B \to A$ is an **inverse mapping** of $f : A \to B$ if $g \circ f = id_A$ and $f \circ g = id_B$; in other words, the inverse function of a function simply "undoes" the function. A map is said to be **invertible** if it has an inverse. We usually write f^{-1} for the inverse of f.

Example 1.16. The function $f(x) = x^3$ has inverse $f^{-1}(x) = \sqrt[3]{x}$ by Example 1.12.

Example 1.17. The natural logarithm and the exponential functions, $f(x) = \ln x$ and $f^{-1}(x) = e^x$, are inverses of each other provided that we are careful about choosing domains. Observe that
$$f(f^{-1}(x)) = f(e^x) = \ln e^x = x$$

and
$$f^{-1}(f(x)) = f^{-1}(\ln x) = e^{\ln x} = x$$

whenever composition makes sense.

Example 1.18. Suppose that

$$A = \begin{pmatrix} 3 & 1 \\ 5 & 2 \end{pmatrix}.$$

Then A defines a map from $\mathbb{R}^2$ to $\mathbb{R}^2$ by

$$T_A(x, y) = (3x + y, 5x + 2y).$$

We can find an inverse map of T_A by simply inverting the matrix A; that is, $T_A^{-1} = T_{A^{-1}}$. In this example,

$$A^{-1} = \begin{pmatrix} 2 & -1 \\ -5 & 3 \end{pmatrix};$$

hence, the inverse map is given by

$$T_A^{-1}(x, y) = (2x - y, -5x + 3y).$$

It is easy to check that

$$T_A^{-1} \circ T_A(x, y) = T_A \circ T_A^{-1}(x, y) = (x, y).$$

Not every map has an inverse. If we consider the map

$$T_B(x, y) = (3x, 0)$$

given by the matrix

$$B = \begin{pmatrix} 3 & 0 \\ 0 & 0 \end{pmatrix},$$

then an inverse map would have to be of the form

$$T_B^{-1}(x, y) = (ax + by, cx + dy)$$

and

$$(x, y) = T \circ T_B^{-1}(x, y) = (3ax + 3by, 0)$$

for all x and y. Clearly this is impossible because y might not be 0.

Example 1.19. Given the permutation

$$\pi = \begin{pmatrix} 1 & 2 & 3 \\ 2 & 3 & 1 \end{pmatrix}$$

on $S = \{1, 2, 3\}$, it is easy to see that the permutation defined by

$$\pi^{-1} = \begin{pmatrix} 1 & 2 & 3 \\ 3 & 1 & 2 \end{pmatrix}$$

is the inverse of π. In fact, any bijective mapping possesses an inverse, as we will see in the next theorem.

Theorem 1.20. *A mapping is invertible if and only if it is both one-to-one and onto.*

PROOF.   Suppose first that $f : A \rightarrow B$ is invertible with inverse $g : B \rightarrow A$. Then $g \circ f = id_A$ is the identity map; that is, $g(f(a)) = a$. If $a_1, a_2 \in A$ with $f(a_1) = f(a_2)$, then $a_1 = g(f(a_1)) = g(f(a_2)) = a_2$. Consequently, f is one-to-one. Now suppose that $b \in B$. To show that f is onto, it is necessary to find an $a \in A$ such that $f(a) = b$, but $f(g(b)) = b$ with $g(b) \in A$. Let $a = g(b)$.

Conversely, let f be bijective and let $b \in B$. Since f is onto, there exists an $a \in A$ such that $f(a) = b$. Because f is one-to-one, a must be unique. Define g by letting $g(b) = a$. We have now constructed the inverse of f. $\square$

Equivalence Relations and Partitions

A fundamental notion in mathematics is that of equality. We can generalize equality with equivalence relations and equivalence classes. An ***equivalence relation*** on a set X is a relation $R \subset X \times X$ such that

- $(x, x) \in R$ for all $x \in X$ (***reflexive property***);

- $(x, y) \in R$ implies $(y, x) \in R$ (***symmetric property***);

- (x, y) and $(y, z) \in R$ imply $(x, z) \in R$ (***transitive property***).

Given an equivalence relation R on a set X, we usually write $x \sim y$ instead of $(x, y) \in R$. If the equivalence relation already has an associated notation such as $=$, $\equiv$, or $\cong$, we will use that notation.

Example 1.21. Let p, q, r, and s be integers, where q and s are nonzero. Define $p/q \sim r/s$ if $ps = qr$. Clearly $\sim$ is reflexive and symmetric. To show that it is also transitive, suppose that $p/q \sim r/s$ and $r/s \sim t/u$, with q, s, and u all nonzero. Then $ps = qr$ and $ru = st$. Therefore,

$$psu = qru = qst.$$

Since $s \neq 0$, $pu = qt$. Consequently, $p/q \sim t/u$.

Example 1.22. Suppose that f and g are differentiable functions on $\mathbb{R}$. We can define an equivalence relation on such functions by letting $f(x) \sim g(x)$ if $f'(x) = g'(x)$. It is clear that $\sim$ is both reflexive and symmetric. To demonstrate transitivity, suppose that $f(x) \sim g(x)$ and $g(x) \sim h(x)$. From calculus we know that $f(x) - g(x) = c_1$ and $g(x) - h(x) = c_2$, where c_1 and c_2 are both constants. Hence,

$$f(x) - h(x) = (f(x) - g(x)) + (g(x) - h(x)) = c_1 - c_2$$

and $f'(x) - h'(x) = 0$. Therefore, $f(x) \sim h(x)$.

Example 1.23. For (x_1, y_1) and (x_2, y_2) in $\mathbb{R}^2$, define $(x_1, y_1) \sim (x_2, y_2)$ if $x_1^2 + y_1^2 = x_2^2 + y_2^2$. Then $\sim$ is an equivalence relation on $\mathbb{R}^2$.

Example 1.24. Let A and B be 2×2 matrices with entries in the real numbers. We can define an equivalence relation on the set of 2×2 matrices, by saying $A \sim B$ if there exists an invertible matrix P such that $PAP^{-1} = B$. For example, if

$$A = \begin{pmatrix} 1 & 2 \\ -1 & 1 \end{pmatrix} \quad \text{and} \quad B = \begin{pmatrix} -18 & 33 \\ -11 & 20 \end{pmatrix},$$

then $A \sim B$ since $PAP^{-1} = B$ for

$$P = \begin{pmatrix} 2 & 5 \\ 1 & 3 \end{pmatrix}.$$

Let I be the 2×2 identity matrix; that is,

$$I = \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}.$$

Then $IAI^{-1} = IAI = A$; therefore, the relation is reflexive. To show symmetry, suppose that $A \sim B$. Then there exists an invertible matrix P such that $PAP^{-1} = B$. So

$$A = P^{-1}BP = P^{-1}B(P^{-1})^{-1}.$$

Finally, suppose that $A \sim B$ and $B \sim C$. Then there exist invertible matrices P and Q such that $PAP^{-1} = B$ and $QBQ^{-1} = C$. Since

$$C = QBQ^{-1} = QPAP^{-1}Q^{-1} = (QP)A(QP)^{-1},$$

the relation is transitive. Two matrices that are equivalent in this manner are said to be **similar**.

A **partition** $\mathcal{P}$ of a set X is a collection of nonempty sets $X_1, X_2, \ldots$ such that $X_i \cap X_j = \emptyset$ for $i \neq j$ and $\bigcup_k X_k = X$. Let $\sim$ be an equivalence relation on a set X and let $x \in X$. Then $[x] = \{y \in X : y \sim x\}$ is called the **equivalence class** of x. We will see that an equivalence relation gives rise to a partition via equivalence classes. Also, whenever a partition of a set exists, there is some natural underlying equivalence relation, as the following theorem demonstrates.

Theorem 1.25. *Given an equivalence relation $\sim$ on a set X, the equivalence classes of X form a partition of X. Conversely, if $\mathcal{P} = \{X_i\}$ is a partition of a set X, then there is an equivalence relation on X with equivalence classes X_i.*

PROOF.   Suppose there exists an equivalence relation $\sim$ on the set X. For any $x \in X$, the reflexive property shows that $x \in [x]$ and so $[x]$ is nonempty. Clearly $X = \bigcup_{x \in X}[x]$. Now let $x, y \in X$. We need to show that either $[x] = [y]$ or $[x] \cap [y] = \emptyset$. Suppose that the intersection of $[x]$ and $[y]$ is not empty and that $z \in [x] \cap [y]$. Then $z \sim x$ and $z \sim y$. By symmetry and transitivity $x \sim y$; hence, $[x] \subset [y]$. Similarly, $[y] \subset [x]$ and so $[x] = [y]$. Therefore, any two equivalence classes are either disjoint or exactly the same.

Conversely, suppose that $\mathcal{P} = \{X_i\}$ is a partition of a set X. Let two elements be equivalent if they are in the same partition. Clearly, the relation is reflexive. If x is in the same partition as y, then y is in the same partition as x, so $x \sim y$ implies $y \sim x$. Finally, if x is in the same partition as y and y is in the same partition as z, then x must be in the same partition as z, and transitivity holds. $\square$

Corollary 1.26. *Two equivalence classes of an equivalence relation are either disjoint or equal.*

Let us examine some of the partitions given by the equivalence classes in the last set of examples.

Example 1.27. In the equivalence relation in Example 1.21, two pairs of integers, (p, q) and (r, s), are in the same equivalence class when they reduce to the same fraction in its lowest terms.

Example 1.28. In the equivalence relation in Example 1.22, two functions $f(x)$ and $g(x)$ are in the same partition when they differ by a constant.

Example 1.29. We defined an equivalence class on $\mathbb{R}^2$ by $(x_1, y_1) \sim (x_2, y_2)$ if $x_1^2 + y_1^2 = x_2^2 + y_2^2$. Two pairs of real numbers are in the same partition when they lie on the same circle about the origin.

Example 1.30. Let r and s be two integers and suppose that $n \in \mathbb{N}$. We say that r is **congruent** to s **modulo** n, or r is congruent to s mod n, if $r - s$ is evenly divisible by n; that is, $r - s = nk$ for some $k \in \mathbb{Z}$. In this case we write $r \equiv s \pmod{n}$. For example, $41 \equiv 17 \pmod{8}$ since $41 - 17 = 24$ is divisible by 8. We claim that congruence modulo n forms an equivalence relation of $\mathbb{Z}$. Certainly any integer r is equivalent to itself since $r - r = 0$ is divisible by n. We will now show that the relation is symmetric. If $r \equiv s$

(mod n), then $r - s = -(s - r)$ is divisible by n. So $s - r$ is divisible by n and $s \equiv r$ (mod n). Now suppose that $r \equiv s$ (mod n) and $s \equiv t$ (mod n). Then there exist integers k and l such that $r - s = kn$ and $s - t = ln$. To show transitivity, it is necessary to prove that $r - t$ is divisible by n. However,

$$r - t = r - s + s - t = kn + ln = (k + l)n,$$

and so $r - t$ is divisible by n.

   If we consider the equivalence relation established by the integers modulo 3, then

$$[0] = \{\ldots, -3, 0, 3, 6, \ldots\},$$
$$[1] = \{\ldots, -2, 1, 4, 7, \ldots\},$$
$$[2] = \{\ldots, -1, 2, 5, 8, \ldots\}.$$

Notice that $[0] \cup [1] \cup [2] = \mathbb{Z}$ and also that the sets are disjoint. The sets $[0]$, $[1]$, and $[2]$ form a partition of the integers.

   The integers modulo n are a very important example in the study of abstract algebra and will become quite useful in our investigation of various algebraic structures such as groups and rings. In our discussion of the integers modulo n we have actually assumed a result known as the division algorithm, which will be stated and proved in Chapter 2.

1.3 Exercises

1.   Suppose that

$$A = \{x : x \in \mathbb{N} \text{ and } x \text{ is even}\},$$
$$B = \{x : x \in \mathbb{N} \text{ and } x \text{ is prime}\},$$
$$C = \{x : x \in \mathbb{N} \text{ and } x \text{ is a multiple of } 5\}.$$

Describe each of the following sets.

   (a) $A \cap B$ (c) $A \cup B$

   (b) $B \cap C$ (d) $A \cap (B \cup C)$

2.   If $A = \{a, b, c\}$, $B = \{1, 2, 3\}$, $C = \{x\}$, and $D = \emptyset$, list all of the elements in each of the following sets.

   (a) $A \times B$ (c) $A \times B \times C$

   (b) $B \times A$ (d) $A \times D$

3.   Find an example of two nonempty sets A and B for which $A \times B = B \times A$ is true.

4.   Prove $A \cup \emptyset = A$ and $A \cap \emptyset = \emptyset$.

5.   Prove $A \cup B = B \cup A$ and $A \cap B = B \cap A$.

6.   Prove $A \cup (B \cap C) = (A \cup B) \cap (A \cup C)$.

7.   Prove $A \cap (B \cup C) = (A \cap B) \cup (A \cap C)$.

8.   Prove $A \subset B$ if and only if $A \cap B = A$.

9.   Prove $(A \cap B)' = A' \cup B'$.

10.   Prove $A \cup B = (A \cap B) \cup (A \setminus B) \cup (B \setminus A)$.

11.   Prove $(A \cup B) \times C = (A \times C) \cup (B \times C)$.

12.   Prove $(A \cap B) \setminus B = \emptyset$.

13.   Prove $(A \cup B) \setminus B = A \setminus B$.

14.   Prove $A \setminus (B \cup C) = (A \setminus B) \cap (A \setminus C)$.

15.   Prove $A \cap (B \setminus C) = (A \cap B) \setminus (A \cap C)$.

16.   Prove $(A \setminus B) \cup (B \setminus A) = (A \cup B) \setminus (A \cap B)$.

17.   Which of the following relations $f : \mathbb{Q} \to \mathbb{Q}$ define a mapping? In each case, supply a reason why f is or is not a mapping.

(a) $f(p/q) = \dfrac{p+1}{p-2}$

(b) $f(p/q) = \dfrac{3p}{3q}$

(c) $f(p/q) = \dfrac{p+q}{q^2}$

(d) $f(p/q) = \dfrac{3p^2}{7q^2} - \dfrac{p}{q}$

18.   Determine which of the following functions are one-to-one and which are onto. If the function is not onto, determine its range.

(a) $f : \mathbb{R} \to \mathbb{R}$ defined by $f(x) = e^x$

(b) $f : \mathbb{Z} \to \mathbb{Z}$ defined by $f(n) = n^2 + 3$

(c) $f : \mathbb{R} \to \mathbb{R}$ defined by $f(x) = \sin x$

(d) $f : \mathbb{Z} \to \mathbb{Z}$ defined by $f(x) = x^2$

19.   Let $f : A \to B$ and $g : B \to C$ be invertible mappings; that is, mappings such that f^{-1} and g^{-1} exist. Show that $(g \circ f)^{-1} = f^{-1} \circ g^{-1}$.

20.

(a) Define a function $f : \mathbb{N} \to \mathbb{N}$ that is one-to-one but not onto.

(b) Define a function $f : \mathbb{N} \to \mathbb{N}$ that is onto but not one-to-one.

21.   Prove the relation defined on $\mathbb{R}^2$ by $(x_1, y_1) \sim (x_2, y_2)$ if $x_1^2 + y_1^2 = x_2^2 + y_2^2$ is an equivalence relation.

22.   Let $f : A \to B$ and $g : B \to C$ be maps.

(a) If f and g are both one-to-one functions, show that $g \circ f$ is one-to-one.

(b) If $g \circ f$ is onto, show that g is onto.

(c) If $g \circ f$ is one-to-one, show that f is one-to-one.

(d) If $g \circ f$ is one-to-one and f is onto, show that g is one-to-one.

(e) If $g \circ f$ is onto and g is one-to-one, show that f is onto.

23.   Define a function on the real numbers by

$$f(x) = \frac{x+1}{x-1}.$$

What are the domain and range of f? What is the inverse of f? Compute $f \circ f^{-1}$ and $f^{-1} \circ f$.

24.   Let $f : X \to Y$ be a map with $A_1, A_2 \subset X$ and $B_1, B_2 \subset Y$.

(a) Prove $f(A_1 \cup A_2) = f(A_1) \cup f(A_2)$.

(b) Prove $f(A_1 \cap A_2) \subset f(A_1) \cap f(A_2)$. Give an example in which equality fails.

(c) Prove $f^{-1}(B_1 \cup B_2) = f^{-1}(B_1) \cup f^{-1}(B_2)$, where

$$f^{-1}(B) = \{x \in X : f(x) \in B\}.$$

(d) Prove $f^{-1}(B_1 \cap B_2) = f^{-1}(B_1) \cap f^{-1}(B_2)$.

(e) Prove $f^{-1}(Y \setminus B_1) = X \setminus f^{-1}(B_1)$.

25. Determine whether or not the following relations are equivalence relations on the given set. If the relation is an equivalence relation, describe the partition given by it. If the relation is not an equivalence relation, state why it fails to be one.

(a) $x \sim y$ in $\mathbb{R}$ if $x \geq y$ (c) $x \sim y$ in $\mathbb{R}$ if $|x - y| \leq 4$

(b) $m \sim n$ in $\mathbb{Z}$ if $mn > 0$ (d) $m \sim n$ in $\mathbb{Z}$ if $m \equiv n \pmod 6$

26. Define a relation $\sim$ on $\mathbb{R}^2$ by stating that $(a, b) \sim (c, d)$ if and only if $a^2 + b^2 \leq c^2 + d^2$. Show that $\sim$ is reflexive and transitive but not symmetric.

27. Show that an $m \times n$ matrix gives rise to a well-defined map from $\mathbb{R}^n$ to $\mathbb{R}^m$.

28. Find the error in the following argument by providing a counterexample. "The reflexive property is redundant in the axioms for an equivalence relation. If $x \sim y$, then $y \sim x$ by the symmetric property. Using the transitive property, we can deduce that $x \sim x$."

29. **Projective Real Line.** Define a relation on $\mathbb{R}^2 \setminus \{(0,0)\}$ by letting $(x_1, y_1) \sim (x_2, y_2)$ if there exists a nonzero real number λ such that $(x_1, y_1) = (\lambda x_2, \lambda y_2)$. Prove that $\sim$ defines an equivalence relation on $\mathbb{R}^2 \setminus (0,0)$. What are the corresponding equivalence classes? This equivalence relation defines the projective line, denoted by $\mathbb{P}(\mathbb{R})$, which is very important in geometry.

References and Suggested Readings

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1.4   Sage

Sage is a powerful system for studying and exploring many different areas of mathematics. In this textbook, you will study a variety of algebraic structures, such as groups, rings and fields. Sage does an excellent job of implementing many features of these objects as we will see in the chapters ahead. But here and now, in this initial chapter, we will concentrate on a few general ways of getting the most out of working with Sage.

You may use Sage several different ways. It may be used as a command-line program when installed on your own computer. Or it might be a web application such as the SageMathCloud. Our writing will assume that you are reading this as a worksheet within the Sage Notebook (a web browser interface), or this is a section of the entire book presented as web pages, and you are employing the Sage Cell Server via those pages. After the first few chapters the explanations should work equally well for whatever vehicle you use to execute Sage commands.

Executing Sage Commands

Most of your interaction will be by typing commands into a *compute cell.* If you are reading this in the Sage Notebook or as a webpage version of the book, then you will see a compute cell just below this paragraph. Click once inside the compute cell and if you are in the Sage Notebook, you will get a more distinctive border around it, a blinking cursor inside, plus a cute little "evaluate" link below.At the cursor, type `2+2` and then click on the evaluate link. Did a `4` appear below the cell? If so, you have successfully sent a command off for Sage to evaluate and you have received back the (correct) answer.

Here is another compute cell. Try evaluating the command `factorial(300)` here.Hmmmmm. That is quite a big integer! If you see slashes at the end of each line, this means the result is continued onto the next line, since there are 615 total digits in the result.

To make new compute cells in the Sage Notebook (only), hover your mouse just above another compute cell, or just below some output from a compute cell. When you see a skinny blue bar across the width of your worksheet, click and you will open up a new compute cell, ready for input. Note that your worksheet will remember any calculations you make, in the order you make them, no matter where you put the cells, so it is best to stay organized and add new cells at the bottom.

Try placing your cursor just below the monstrous value of 300! that you have. Click on the blue bar and try another factorial computation in the new compute cell.

Each compute cell will show output due to only the very last command in the cell. Try to predict the following output before evaluating the cell.

```
a = 10
b = 6
b = b - 10
```

```
a = a + 20
a
```

30

The following compute cell will not print anything since the one command does not create output. But it will have an effect, as you can see when you execute the subsequent cell. Notice how this uses the value of b from above. Execute this compute cell *once*. Exactly once. Even if it *appears* to do nothing. If you execute the cell twice, your credit card may be charged twice.

```
b = b + 50
```

Now execute this cell, which will produce some output.

```
b + 20
```

66

So b came into existence as 6. We subtracted 10 immediately afterward. Then a subsequent cell added 50. This assumes you executed this cell *exactly* once! In the last cell we create b+20 (but do not save it) and it is this value (66) that is output, while b is still 46.

You can combine several commands on one line with a semi-colon. This is a great way to get multiple outputs from a compute cell. The syntax for building a matrix should be somewhat obvious when you see the output, but if not, it is not particularly important to understand now.

```
A = matrix([[3, 1], [5,2]]); A
```

```
[3 1]
[5 2]
```

```
print(A); print(); print(A.inverse())
```

```
[3 1]
[5 2]
<BLANKLINE>
[ 2 -1]
[-5  3]
```

Immediate Help

Some commands in Sage are "functions," an example is `factorial()` above. Other commands are "methods" of an object and are like characteristics of objects, an example is `.inverse()` as a method of a matrix. Once you know how to create an object (such as a matrix), then it is easy to see all the available methods. Write the name of the object, place a period ("dot") and hit the TAB key. If you have A defined from above, then the compute cell below is ready to go, click into it and then hit TAB (not "evaluate"!). You should get a long list of possible methods.

```
A.
```

To get some help on how to use a method with an object, write its name after a dot (with no parentheses) and then use a question-mark and hit TAB. (Hit the escape key "ESC" to remove the list, or click on the text for a method.)

```
A.inverse?
```

With one more question-mark and a TAB you can see the actual computer instructions that were programmed into Sage to make the method work, once you scoll down past the documentation delimited by the triple quotes ("""):

```
A.inverse??
```

It is worthwhile to see what Sage does when there is an error. You will probably see a lot of these at first, and initially they will be a bit intimidating. But with time, you will learn how to use them effectively and you will also become more proficient with Sage and see them less often. Execute the compute cell below, it asks for the inverse of a matrix that has no inverse. Then reread the commentary.

```
B = matrix([[2, 20], [5, 50]])
B.inverse()
```

```
Traceback (most recent call last):
...
ZeroDivisionError: Matrix is singular
```

Click just to the left of the error message to expand it fully (another click hides it totally, and a third click brings back the abbreviated form). Read the bottom of an error message first, it is your best explanation. Here a `ZeroDivisionError` is not 100% accurate, but is close. The matrix is not invertible, not dissimilar to how we cannot divide scalars by zero. The remainder of the message begins at the top showing were the error first happened in your code and then the various places where intermediate functions were called, until the actual piece of Sage where the problem occurred. Sometimes this information will give you some clues, sometimes it is totally undecipherable. So do not let it scare you if it seems mysterious, but do remember to always read the last line first, then go back and read the first few lines for something that looks like your code.

Annotating Your Work

It is easy to comment on your work when you use the Sage Notebook. (The following only applies if you are reading this within a Sage Notebook. If you are not, then perhaps you can go open up a worksheet in the Sage Notebook and experiment there.) You can open up a small word-processor by hovering your mouse until you get a skinny blue bar again, but now when you click, also hold the SHIFT key at the same time. Experiment with fonts, colors, bullet lists, etc and then click the "Save changes" button to exit. Double-click on your text if you need to go back and edit it later.

Open the word-processor again to create a new bit of text (maybe next to the empty compute cell just below). Type all of the following *exactly,*

<div align="center">Pythagorean Theorem: $c^2=a^2+b^2$</div>

and save your changes. The symbols between the dollar signs are written according to the mathematical typesetting language known as TEX — cruise the internet to learn more about this very popular tool. (Well, it is extremely popular among mathematicians and physical scientists.)

Lists

Much of our interaction with sets will be through Sage lists. These are not really sets — they allow duplicates, and order matters. But they are so close to sets, and so easy and powerful to use that we will use them regularly. We will use a fun made-up list for practice, the

quote marks mean the items are just text, with no special mathematical meaning. Execute these compute cells as we work through them.

```
zoo = ['snake', 'parrot', 'elephant', 'baboon', 'beetle']
zoo
```

```
['snake', 'parrot', 'elephant', 'baboon', 'beetle']
```

So the square brackets define the boundaries of our list, commas separate items, and we can give the list a name. To work with just one element of the list, we use the name and a pair of brackets with an index. Notice that lists have indices that *begin counting at zero*. This will seem odd at first and will seem very natural later.

```
zoo[2]
```

```
'elephant'
```

We can add a new creature to the zoo, it is joined up at the far right end.

```
zoo.append('ostrich'); zoo
```

```
['snake', 'parrot', 'elephant', 'baboon', 'beetle', 'ostrich']
```

We can remove a creature.

```
zoo.remove('parrot')
zoo
```

```
['snake', 'elephant', 'baboon', 'beetle', 'ostrich']
```

We can extract a sublist. Here we start with element 1 (the elephant) and go all the way up to, *but not including*, element 3 (the beetle). Again a bit odd, but it will feel natural later. For now, notice that we are extracting two elements of the lists, exactly $3 - 1 = 2$ elements.

```
mammals = zoo[1:3]
mammals
```

```
['elephant', 'baboon']
```

Often we will want to see if two lists are equal. To do that we will need to sort a list first. A function creates a new, sorted list, leaving the original alone. So we need to save the new one with a new name.

```
newzoo = sorted(zoo)
newzoo
```

```
['baboon', 'beetle', 'elephant', 'ostrich', 'snake']
```

```
zoo.sort()
zoo
```

```
['baboon', 'beetle', 'elephant', 'ostrich', 'snake']
```

Notice that if you run this last compute cell your zoo has changed and some commands above will not necessarily execute the same way. If you want to experiment, go all the way back to the first creation of the zoo and start executing cells again from there with a fresh zoo.

A construction called a ***list comprehension*** is especially powerful, especially since it almost exactly mirrors notation we use to describe sets. Suppose we want to form the plural of the names of the creatures in our zoo. We build a new list, based on all of the elements of our old list.

```
plurality_zoo = [animal+'s' for animal in zoo]
plurality_zoo
```

```
['baboons', 'beetles', 'elephants', 'ostrichs', 'snakes']
```

Almost like it says: we add an "s" to each animal name, for each animal in the zoo, and place them in a new list. Perfect. (Except for getting the plural of "ostrich" wrong.)

Lists of Integers

One final type of list, with numbers this time. The `srange()` function will create lists of integers. (The "s" in the name stands for "Sage" and so will produce integers that Sage understands best. Many early difficulties with Sage and group theory can be alleviated by using only this command to create lists of integers.) In its simplest form an invocation like `srange(12)` will create a list of 12 integers, *starting at zero* and working up to, *but not including*, 12. Does this sound familiar?

```
dozen = srange(12); dozen
```

```
[0, 1, 2, 3, 4, 5, 6, 7, 8, 9, 10, 11]
```

Here are two other forms, that you should be able to understand by studying the examples.

```
teens = srange(13, 20); teens
```

```
[13, 14, 15, 16, 17, 18, 19]
```

```
decades = srange(1900, 2000, 10); decades
```

```
[1900, 1910, 1920, 1930, 1940, 1950, 1960, 1970, 1980, 1990]
```

Saving and Sharing Your Work

There is a "Save" button in the upper-right corner of the Sage Notebook. This will save a current copy of your worksheet that you can retrieve your work from within your notebook again later, though you have to re-execute all the cells when you re-open the worksheet.

There is also a "File" drop-down list, on the left, just above your very top compute cell (not be confused with your browser's File menu item!). You will see a choice here labeled "Save worksheet to a file..." When you do this, you are creating a copy of your worksheet in the `sws` format (short for "Sage WorkSheet"). You can email this file, or post it on a website, for other Sage users and they can use the "Upload" link on the homepage of their notebook to incorporate a copy of your worksheet into their notebook.

There are other ways to share worksheets that you can experiment with, but this gives you one way to share any worksheet with anybody almost anywhere.

We have covered a lot here in this section, so come back later to pick up tidbits you might have missed. There are also many more features in the Sage Notebook that we have not covered.

1.5   Sage Exercises

1.   This exercise is just about making sure you know how to use Sage. You may be using the Sage Notebook server the online CoCalc service through your web browser. In either event, create a new worksheet. Do some non-trivial computation, maybe a pretty

plot or some gruesome numerical computation to an insane precision. Create an interesting list and experiment with it some. Maybe include some nicely formatted text or TeX using the included mini-word-processor of the Sage Notebook (hover until a blue bar appears between cells and then shift-click) or create commentary in cells within CoCalc using the *magics* %html or %md on a line of their own followed by text in HTML or Markdown syntax (respectively).

Use whatever mechanism your instructor has in place for submitting your work. Or save your worksheet and then trade with a classmate.

$$\boxed{2}$$

The Integers

The integers are the building blocks of mathematics. In this chapter we will investigate the fundamental properties of the integers, including mathematical induction, the division algorithm, and the Fundamental Theorem of Arithmetic.

2.1 Mathematical Induction

Suppose we wish to show that

$$1 + 2 + \cdots + n = \frac{n(n+1)}{2}$$

for any natural number n. This formula is easily verified for small numbers such as $n = 1$, 2, 3, or 4, but it is impossible to verify for all natural numbers on a case-by-case basis. To prove the formula true in general, a more generic method is required.

Suppose we have verified the equation for the first n cases. We will attempt to show that we can generate the formula for the $(n+1)$th case from this knowledge. The formula is true for $n = 1$ since

$$1 = \frac{1(1+1)}{2}.$$

If we have verified the first n cases, then

$$\begin{aligned}
1 + 2 + \cdots + n + (n+1) &= \frac{n(n+1)}{2} + n + 1 \\
&= \frac{n^2 + 3n + 2}{2} \\
&= \frac{(n+1)[(n+1)+1]}{2}.
\end{aligned}$$

This is exactly the formula for the $(n+1)$th case.

This method of proof is known as ***mathematical induction***. Instead of attempting to verify a statement about some subset S of the positive integers $\mathbb{N}$ on a case-by-case basis, an impossible task if S is an infinite set, we give a specific proof for the smallest integer being considered, followed by a generic argument showing that if the statement holds for a given case, then it must also hold for the next case in the sequence. We summarize mathematical induction in the following axiom.

Principle 2.1 First Principle of Mathematical Induction. *Let $S(n)$ be a statement about integers for $n \in \mathbb{N}$ and suppose $S(n_0)$ is true for some integer n_0. If for all integers k with $k \geq n_0$, $S(k)$ implies that $S(k+1)$ is true, then $S(n)$ is true for all integers n greater than or equal to n_0.*

Example 2.2. For all integers $n \geq 3$, $2^n > n + 4$. Since

$$8 = 2^3 > 3 + 4 = 7,$$

the statement is true for $n_0 = 3$. Assume that $2^k > k + 4$ for $k \geq 3$. Then $2^{k+1} = 2 \cdot 2^k > 2(k + 4)$. But

$$2(k + 4) = 2k + 8 > k + 5 = (k + 1) + 4$$

since k is positive. Hence, by induction, the statement holds for all integers $n \geq 3$.

Example 2.3. Every integer $10^{n+1} + 3 \cdot 10^n + 5$ is divisible by 9 for $n \in \mathbb{N}$. For $n = 1$,

$$10^{1+1} + 3 \cdot 10 + 5 = 135 = 9 \cdot 15$$

is divisible by 9. Suppose that $10^{k+1} + 3 \cdot 10^k + 5$ is divisible by 9 for $k \geq 1$. Then

$$10^{(k+1)+1} + 3 \cdot 10^{k+1} + 5 = 10^{k+2} + 3 \cdot 10^{k+1} + 50 - 45$$
$$= 10(10^{k+1} + 3 \cdot 10^k + 5) - 45$$

is divisible by 9.

Example 2.4. We will prove the binomial theorem using mathematical induction; that is,

$$(a + b)^n = \sum_{k=0}^{n} \binom{n}{k} a^k b^{n-k},$$

where a and b are real numbers, $n \in \mathbb{N}$, and

$$\binom{n}{k} = \frac{n!}{k!(n-k)!}$$

is the binomial coefficient. We first show that

$$\binom{n+1}{k} = \binom{n}{k} + \binom{n}{k-1}.$$

This result follows from

$$\binom{n}{k} + \binom{n}{k-1} = \frac{n!}{k!(n-k)!} + \frac{n!}{(k-1)!(n-k+1)!}$$
$$= \frac{(n+1)!}{k!(n+1-k)!}$$
$$= \binom{n+1}{k}.$$

If $n = 1$, the binomial theorem is easy to verify. Now assume that the result is true for n greater than or equal to 1. Then

$$(a + b)^{n+1} = (a + b)(a + b)^n$$
$$= (a + b)\left(\sum_{k=0}^{n} \binom{n}{k} a^k b^{n-k} \right)$$
$$= \sum_{k=0}^{n} \binom{n}{k} a^{k+1} b^{n-k} + \sum_{k=0}^{n} \binom{n}{k} a^k b^{n+1-k}$$

$$= a^{n+1} + \sum_{k=1}^{n} \binom{n}{k-1} a^k b^{n+1-k} + \sum_{k=1}^{n} \binom{n}{k} a^k b^{n+1-k} + b^{n+1}$$

$$= a^{n+1} + \sum_{k=1}^{n} \left[\binom{n}{k-1} + \binom{n}{k} \right] a^k b^{n+1-k} + b^{n+1}$$

$$= \sum_{k=0}^{n+1} \binom{n+1}{k} a^k b^{n+1-k}.$$

We have an equivalent statement of the Principle of Mathematical Induction that is often very useful.

Principle 2.5 Second Principle of Mathematical Induction. *Let $S(n)$ be a statement about integers for $n \in \mathbb{N}$ and suppose $S(n_0)$ is true for some integer n_0. If $S(n_0), S(n_0 + 1), \ldots, S(k)$ imply that $S(k+1)$ for $k \geq n_0$, then the statement $S(n)$ is true for all integers $n \geq n_0$.*

A nonempty subset S of $\mathbb{Z}$ is **well-ordered** if S contains a least element. Notice that the set $\mathbb{Z}$ is not well-ordered since it does not contain a smallest element. However, the natural numbers are well-ordered.

Principle 2.6 Principle of Well-Ordering. *Every nonempty subset of the natural numbers is well-ordered.*

The Principle of Well-Ordering is equivalent to the Principle of Mathematical Induction.

Lemma 2.7. *The Principle of Mathematical Induction implies that 1 is the least positive natural number.*

PROOF.   Let $S = \{n \in \mathbb{N} : n \geq 1\}$. Then $1 \in S$. Assume that $n \in S$. Since $0 < 1$, it must be the case that $n = n + 0 < n + 1$. Therefore, $1 \leq n < n + 1$. Consequently, if $n \in S$, then $n + 1$ must also be in S, and by the Principle of Mathematical Induction, and $S = \mathbb{N}$. $\square$

Theorem 2.8. *The Principle of Mathematical Induction implies the Principle of Well-Ordering. That is, every nonempty subset of $\mathbb{N}$ contains a least element.*

PROOF.   We must show that if S is a nonempty subset of the natural numbers, then S contains a least element. If S contains 1, then the theorem is true by Lemma 2.7. Assume that if S contains an integer k such that $1 \leq k \leq n$, then S contains a least element. We will show that if a set S contains an integer less than or equal to $n + 1$, then S has a least element. If S does not contain an integer less than $n + 1$, then $n + 1$ is the smallest integer in S. Otherwise, since S is nonempty, S must contain an integer less than or equal to n. In this case, by induction, S contains a least element. $\square$

Induction can also be very useful in formulating definitions. For instance, there are two ways to define $n!$, the factorial of a positive integer n.

- The *explicit* definition: $n! = 1 \cdot 2 \cdot 3 \cdots (n-1) \cdot n$.

- The *inductive* or *recursive* definition: $1! = 1$ and $n! = n(n-1)!$ for $n > 1$.

Every good mathematician or computer scientist knows that looking at problems recursively, as opposed to explicitly, often results in better understanding of complex issues.

2.2   The Division Algorithm

An application of the Principle of Well-Ordering that we will use often is the division algorithm.

Theorem 2.9 Division Algorithm. *Let a and b be integers, with $b > 0$. Then there exist unique integers q and r such that*

$$a = bq + r$$

where $0 \leq r < b$.

PROOF.   This is a perfect example of the existence-and-uniqueness type of proof. We must first prove that the numbers q and r actually exist. Then we must show that if q' and r' are two other such numbers, then $q = q'$ and $r = r'$.

   Existence of q and r. Let

$$S = \{a - bk : k \in \mathbb{Z} \text{ and } a - bk \geq 0\}.$$

If $0 \in S$, then b divides a, and we can let $q = a/b$ and $r = 0$. If $0 \notin S$, we can use the Well-Ordering Principle. We must first show that S is nonempty. If $a > 0$, then $a - b \cdot 0 \in S$. If $a < 0$, then $a - b(2a) = a(1 - 2b) \in S$. In either case $S \neq \emptyset$. By the Well-Ordering Principle, S must have a smallest member, say $r = a - bq$. Therefore, $a = bq + r$, $r \geq 0$. We now show that $r < b$. Suppose that $r > b$. Then

$$a - b(q + 1) = a - bq - b = r - b > 0.$$

In this case we would have $a - b(q + 1)$ in the set S. But then $a - b(q + 1) < a - bq$, which would contradict the fact that $r = a - bq$ is the smallest member of S. So $r \leq b$. Since $0 \notin S$, $r \neq b$ and so $r < b$.

   Uniqueness of q and r. Suppose there exist integers r, r', q, and q' such that

$$a = bq + r, 0 \leq r < b \quad \text{and} \quad a = bq' + r', 0 \leq r' < b.$$

Then $bq + r = bq' + r'$. Assume that $r' \geq r$. From the last equation we have $b(q - q') = r' - r$; therefore, b must divide $r' - r$ and $0 \leq r' - r \leq r' < b$. This is possible only if $r' - r = 0$. Hence, $r = r'$ and $q = q'$.   □

   Let a and b be integers. If $b = ak$ for some integer k, we write $a \mid b$. An integer d is called a ***common divisor*** of a and b if $d \mid a$ and $d \mid b$. The ***greatest common divisor*** of integers a and b is a positive integer d such that d is a common divisor of a and b and if d' is any other common divisor of a and b, then $d' \mid d$. We write $d = \gcd(a, b)$; for example, $\gcd(24, 36) = 12$ and $\gcd(120, 102) = 6$. We say that two integers a and b are ***relatively prime*** if $\gcd(a, b) = 1$.

Theorem 2.10. *Let a and b be nonzero integers. Then there exist integers r and s such that*

$$\gcd(a, b) = ar + bs.$$

Furthermore, the greatest common divisor of a and b is unique.

PROOF.   Let

$$S = \{am + bn : m, n \in \mathbb{Z} \text{ and } am + bn > 0\}.$$

Clearly, the set S is nonempty; hence, by the Well-Ordering Principle S must have a smallest member, say $d = ar + bs$. We claim that $d = \gcd(a, b)$. Write $a = dq + r'$ where $0 \leq r' < d$. If $r' > 0$, then

$$r' = a - dq$$

$$= a - (ar + bs)q$$
$$= a - arq - bsq$$
$$= a(1 - rq) + b(-sq),$$

which is in S. But this would contradict the fact that d is the smallest member of S. Hence, $r' = 0$ and d divides a. A similar argument shows that d divides b. Therefore, d is a common divisor of a and b.

Suppose that d' is another common divisor of a and b, and we want to show that $d' \mid d$. If we let $a = d'h$ and $b = d'k$, then

$$d = ar + bs = d'hr + d'ks = d'(hr + ks).$$

So d' must divide d. Hence, d must be the unique greatest common divisor of a and b. $\square$

Corollary 2.11. *Let a and b be two integers that are relatively prime. Then there exist integers r and s such that $ar + bs = 1$.*

The Euclidean Algorithm

Among other things, Theorem 2.10 allows us to compute the greatest common divisor of two integers.

Example 2.12. Let us compute the greatest common divisor of 945 and 2415. First observe that

$$2415 = 945 \cdot 2 + 525$$
$$945 = 525 \cdot 1 + 420$$
$$525 = 420 \cdot 1 + 105$$
$$420 = 105 \cdot 4 + 0.$$

Reversing our steps, 105 divides 420, 105 divides 525, 105 divides 945, and 105 divides 2415. Hence, 105 divides both 945 and 2415. If d were another common divisor of 945 and 2415, then d would also have to divide 105. Therefore, $\gcd(945, 2415) = 105$.

If we work backward through the above sequence of equations, we can also obtain numbers r and s such that $945r + 2415s = 105$. Observe that

$$105 = 525 + (-1) \cdot 420$$
$$= 525 + (-1) \cdot [945 + (-1) \cdot 525]$$
$$= 2 \cdot 525 + (-1) \cdot 945$$
$$= 2 \cdot [2415 + (-2) \cdot 945] + (-1) \cdot 945$$
$$= 2 \cdot 2415 + (-5) \cdot 945.$$

So $r = -5$ and $s = 2$. Notice that r and s are not unique, since $r = 41$ and $s = -16$ would also work.

To compute $\gcd(a, b) = d$, we are using repeated divisions to obtain a decreasing sequence of positive integers $r_1 > r_2 > \cdots > r_n = d$; that is,

$$b = aq_1 + r_1$$
$$a = r_1 q_2 + r_2$$
$$r_1 = r_2 q_3 + r_3$$

$$\vdots$$
$$r_{n-2} = r_{n-1}q_n + r_n$$
$$r_{n-1} = r_n q_{n+1}.$$

To find r and s such that $ar + bs = d$, we begin with this last equation and substitute results obtained from the previous equations:

$$d = r_n$$
$$= r_{n-2} - r_{n-1}q_n$$
$$= r_{n-2} - q_n(r_{n-3} - q_{n-1}r_{n-2})$$
$$= -q_n r_{n-3} + (1 + q_n q_{n-1})r_{n-2}$$
$$\vdots$$
$$= ra + sb.$$

The algorithm that we have just used to find the greatest common divisor d of two integers a and b and to write d as the linear combination of a and b is known as the **Euclidean algorithm**.

Prime Numbers

Let p be an integer such that $p > 1$. We say that p is a **prime number**, or simply p is **prime**, if the only positive numbers that divide p are 1 and p itself. An integer $n > 1$ that is not prime is said to be **composite**.

Lemma 2.13 Euclid. *Let a and b be integers and p be a prime number. If $p \mid ab$, then either $p \mid a$ or $p \mid b$.*

PROOF. Suppose that p does not divide a. We must show that $p \mid b$. Since $\gcd(a, p) = 1$, there exist integers r and s such that $ar + ps = 1$. So

$$b = b(ar + ps) = (ab)r + p(bs).$$

Since p divides both ab and itself, p must divide $b = (ab)r + p(bs)$. $\qquad\square$

Theorem 2.14 Euclid. *There exist an infinite number of primes.*

PROOF. We will prove this theorem by contradiction. Suppose that there are only a finite number of primes, say $p_1, p_2, \ldots, p_n$. Let $P = p_1 p_2 \cdots p_n + 1$. Then P must be divisible by some p_i for $1 \le i \le n$. In this case, p_i must divide $P - p_1 p_2 \cdots p_n = 1$, which is a contradiction. Hence, either P is prime or there exists an additional prime number $p \ne p_i$ that divides P. $\qquad\square$

Theorem 2.15 Fundamental Theorem of Arithmetic. *Let n be an integer such that $n > 1$. Then*

$$n = p_1 p_2 \cdots p_k,$$

where $p_1, \ldots, p_k$ are primes (not necessarily distinct). Furthermore, this factorization is unique; that is, if

$$n = q_1 q_2 \cdots q_l,$$

then $k = l$ and the q_i's are just the p_i's rearranged.

PROOF.   *Uniqueness.* To show uniqueness we will use induction on n. The theorem is certainly true for $n = 2$ since in this case n is prime. Now assume that the result holds for all integers m such that $1 \leq m < n$, and

$$n = p_1 p_2 \cdots p_k = q_1 q_2 \cdots q_l,$$

where $p_1 \leq p_2 \leq \cdots \leq p_k$ and $q_1 \leq q_2 \leq \cdots \leq q_l$. By Lemma 2.13, $p_1 \mid q_i$ for some $i = 1, \ldots, l$ and $q_1 \mid p_j$ for some $j = 1, \ldots, k$. Since all of the p_i's and q_i's are prime, $p_1 = q_i$ and $q_1 = p_j$. Hence, $p_1 = q_1$ since $p_1 \leq p_j = q_1 \leq q_i = p_1$. By the induction hypothesis,

$$n' = p_2 \cdots p_k = q_2 \cdots q_l$$

has a unique factorization. Hence, $k = l$ and $q_i = p_i$ for $i = 1, \ldots, k$.

   Existence. To show existence, suppose that there is some integer that cannot be written as the product of primes. Let S be the set of all such numbers. By the Principle of Well-Ordering, S has a smallest number, say a. If the only positive factors of a are a and 1, then a is prime, which is a contradiction. Hence, $a = a_1 a_2$ where $1 < a_1 < a$ and $1 < a_2 < a$. Neither $a_1 \in S$ nor $a_2 \in S$, since a is the smallest element in S. So

$$a_1 = p_1 \cdots p_r$$
$$a_2 = q_1 \cdots q_s.$$

Therefore,
$$a = a_1 a_2 = p_1 \cdots p_r q_1 \cdots q_s.$$

So $a \notin S$, which is a contradiction. $\square$

Historical Note

   Prime numbers were first studied by the ancient Greeks. Two important results from antiquity are Euclid's proof that an infinite number of primes exist and the Sieve of Eratosthenes, a method of computing all of the prime numbers less than a fixed positive integer n. One problem in number theory is to find a function f such that $f(n)$ is prime for each integer n. Pierre Fermat (1601?–1665) conjectured that $2^{2^n} + 1$ was prime for all n, but later it was shown by Leonhard Euler (1707–1783) that

$$2^{2^5} + 1 = 4,294,967,297$$

is a composite number. One of the many unproven conjectures about prime numbers is Goldbach's Conjecture. In a letter to Euler in 1742, Christian Goldbach stated the conjecture that every even integer with the exception of 2 seemed to be the sum of two primes: $4 = 2 + 2$, $6 = 3 + 3$, $8 = 3 + 5$, …. Although the conjecture has been verified for the numbers up through 4×10^{18}, it has yet to be proven in general. Since prime numbers play an important role in public key cryptography, there is currently a great deal of interest in determining whether or not a large number is prime.

2.3   Exercises

1.   Prove that
$$1^2 + 2^2 + \cdots + n^2 = \frac{n(n+1)(2n+1)}{6}$$

for $n \in \mathbb{N}$.

2. Prove that

$$1^3 + 2^3 + \cdots + n^3 = \frac{n^2(n+1)^2}{4}$$

for $n \in \mathbb{N}$.

3. Prove that $n! > 2^n$ for $n \geq 4$.

4. Prove that

$$x + 4x + 7x + \cdots + (3n-2)x = \frac{n(3n-1)x}{2}$$

for $n \in \mathbb{N}$.

5. Prove that $10^{n+1} + 10^n + 1$ is divisible by 3 for $n \in \mathbb{N}$.

6. Prove that $4 \cdot 10^{2n} + 9 \cdot 10^{2n-1} + 5$ is divisible by 99 for $n \in \mathbb{N}$.

7. Show that

$$\sqrt[n]{a_1 a_2 \cdots a_n} \leq \frac{1}{n}\sum_{k=1}^{n} a_k.$$

8. Prove the Leibniz rule for $f^{(n)}(x)$, where $f^{(n)}$ is the nth derivative of f; that is, show that

$$(fg)^{(n)}(x) = \sum_{k=0}^{n} \binom{n}{k} f^{(k)}(x)g^{(n-k)}(x).$$

9. Use induction to prove that $1 + 2 + 2^2 + \cdots + 2^n = 2^{n+1} - 1$ for $n \in \mathbb{N}$.

10. Prove that

$$\frac{1}{2} + \frac{1}{6} + \cdots + \frac{1}{n(n+1)} = \frac{n}{n+1}$$

for $n \in \mathbb{N}$.

11. If x is a nonnegative real number, then show that $(1+x)^n - 1 \geq nx$ for $n = 0, 1, 2, \ldots$.

12. Power Sets. Let X be a set. Define the ***power set*** of X, denoted $\mathcal{P}(X)$, to be the set of all subsets of X. For example,

$$\mathcal{P}(\{a, b\}) = \{\emptyset, \{a\}, \{b\}, \{a, b\}\}.$$

For every positive integer n, show that a set with exactly n elements has a power set with exactly 2^n elements.

13. Prove that the two principles of mathematical induction stated in Section 2.1 are equivalent.

14. Show that the Principle of Well-Ordering for the natural numbers implies that 1 is the smallest natural number. Use this result to show that the Principle of Well-Ordering implies the Principle of Mathematical Induction; that is, show that if $S \subset \mathbb{N}$ such that $1 \in S$ and $n + 1 \in S$ whenever $n \in S$, then $S = \mathbb{N}$.

15. For each of the following pairs of numbers a and b, calculate $\gcd(a, b)$ and find integers r and s such that $\gcd(a, b) = ra + sb$.

(a) 14 and 39

(b) 234 and 165

(c) 1739 and 9923

(d) 471 and 562

(e) 23771 and 19945

(f) -4357 and 3754

16. Let a and b be nonzero integers. If there exist integers r and s such that $ar + bs = 1$, show that a and b are relatively prime.

17. **Fibonacci Numbers.**   The Fibonacci numbers are

$$1, 1, 2, 3, 5, 8, 13, 21, \ldots.$$

We can define them inductively by $f_1 = 1$, $f_2 = 1$, and $f_{n+2} = f_{n+1} + f_n$ for $n \in \mathbb{N}$.

(a) Prove that $f_n < 2^n$.

(b) Prove that $f_{n+1}f_{n-1} = f_n^2 + (-1)^n$, $n \geq 2$.

(c) Prove that $f_n = [(1 + \sqrt{5}\,)^n - (1 - \sqrt{5}\,)^n]/2^n\sqrt{5}$.

(d) Show that $\lim_{n \to \infty} f_n/f_{n+1} = (\sqrt{5} - 1)/2$.

(e) Prove that f_n and f_{n+1} are relatively prime.

18.   Let a and b be integers such that $\gcd(a, b) = 1$. Let r and s be integers such that $ar + bs = 1$. Prove that

$$\gcd(a, s) = \gcd(r, b) = \gcd(r, s) = 1.$$

19.   Let $x, y \in \mathbb{N}$ be relatively prime. If xy is a perfect square, prove that x and y must both be perfect squares.

20.   Using the division algorithm, show that every perfect square is of the form $4k$ or $4k + 1$ for some nonnegative integer k.

21.   Suppose that a, b, r, s are pairwise relatively prime and that

$$a^2 + b^2 = r^2$$
$$a^2 - b^2 = s^2.$$

Prove that a, r, and s are odd and b is even.

22.   Let $n \in \mathbb{N}$. Use the division algorithm to prove that every integer is congruent mod n to precisely one of the integers $0, 1, \ldots, n - 1$. Conclude that if r is an integer, then there is exactly one s in $\mathbb{Z}$ such that $0 \leq s < n$ and $[r] = [s]$. Hence, the integers are indeed partitioned by congruence mod n.

23.   Define the **least common multiple** of two nonzero integers a and b, denoted by $\text{lcm}(a, b)$, to be the nonnegative integer m such that both a and b divide m, and if a and b divide any other integer n, then m also divides n. Prove there exists a unique least common multiple for any two integers a and b.

24.   If $d = \gcd(a, b)$ and $m = \text{lcm}(a, b)$, prove that $dm = |ab|$.

25.   Show that $\text{lcm}(a, b) = ab$ if and only if $\gcd(a, b) = 1$.

26.   Prove that $\gcd(a, c) = \gcd(b, c) = 1$ if and only if $\gcd(ab, c) = 1$ for integers a, b, and c.

27.   Let $a, b, c \in \mathbb{Z}$. Prove that if $\gcd(a, b) = 1$ and $a \mid bc$, then $a \mid c$.

28.   Let $p \geq 2$. Prove that if $2^p - 1$ is prime, then p must also be prime.

29.   Prove that there are an infinite number of primes of the form $6n + 5$.

30.   Prove that there are an infinite number of primes of the form $4n - 1$.

31.   Using the fact that 2 is prime, show that there do not exist integers p and q such that $p^2 = 2q^2$. Demonstrate that therefore $\sqrt{2}$ cannot be a rational number.

2.4 Programming Exercises

1. **The Sieve of Eratosthenes.**   One method of computing all of the prime numbers less than a certain fixed positive integer N is to list all of the numbers n such that $1 < n < N$. Begin by eliminating all of the multiples of 2. Next eliminate all of the multiples of 3. Now eliminate all of the multiples of 5. Notice that 4 has already been crossed out. Continue in this manner, noticing that we do not have to go all the way to N; it suffices to stop at $\sqrt{N}$. Using this method, compute all of the prime numbers less than $N = 250$. We can also use this method to find all of the integers that are relatively prime to an integer N. Simply eliminate the prime factors of N and all of their multiples. Using this method, find all of the numbers that are relatively prime to $N = 120$. Using the Sieve of Eratosthenes, write a program that will compute all of the primes less than an integer N.

2.   Let $\mathbb{N}^0 = \mathbb{N} \cup \{0\}$. Ackermann's function is the function $A : \mathbb{N}^0 \times \mathbb{N}^0 \to \mathbb{N}^0$ defined by the equations

$$A(0, y) = y + 1,$$
$$A(x + 1, 0) = A(x, 1),$$
$$A(x + 1, y + 1) = A(x, A(x + 1, y)).$$

Use this definition to compute $A(3, 1)$. Write a program to evaluate Ackermann's function. Modify the program to count the number of statements executed in the program when Ackermann's function is evaluated. How many statements are executed in the evaluation of $A(4, 1)$? What about $A(5, 1)$?

3.   Write a computer program that will implement the Euclidean algorithm. The program should accept two positive integers a and b as input and should output $\gcd(a, b)$ as well as integers r and s such that

$$\gcd(a, b) = ra + sb.$$

References and Suggested Readings

[1]   Brookshear, J. G. *Theory of Computation: Formal Languages, Automata, and Complexity.* Benjamin/Cummings, Redwood City, CA, 1989. Shows the relationships of the theoretical aspects of computer science to set theory and the integers.

[2]   Hardy, G. H. and Wright, E. M. *An Introduction to the Theory of Numbers.* 6th ed. Oxford University Press, New York, 2008.

[3]   Niven, I. and Zuckerman, H. S. *An Introduction to the Theory of Numbers.* 5th ed. Wiley, New York, 1991.

[4]   Vanden Eynden, C. *Elementary Number Theory.* 2nd ed. Waveland Press, Long Grove IL, 2001.

2.5 Sage

Many properties of the algebraic objects we will study can be determined from properties of associated integers. And Sage has many powerful functions for analyzing integers.

Division Algorithm

The code `a % b` will return the remainder upon division of a by b. In other words, the result is the unique integer r such that (1) $0 \le r < b$, and (2) $a = bq + r$ for some integer q (the quotient), as guaranteed by the Division Algorithm (Theorem 2.9). Then $(a - r)/b$ will equal q. For example,

```
r = 14 % 3
r
```

2

```
q = (14 - r)/3
q
```

4

It is also possible to get both the quotient and remainder at the same time with the `.quo_rem()` method (quotient and remainder).

```
a = 14
b = 3
a.quo_rem(b)
```

(4, 2)

A remainder of zero indicates divisibility. So `(a % b) == 0` will return `True` if b divides a, and will otherwise return `False`.

```
(20 % 5) == 0
```

True

```
(17 % 4) == 0
```

False

The `.divides()` method is another option.

```
c = 5
c.divides(20)
```

True

```
d = 4
d.divides(17)
```

False

Greatest Common Divisor

The greatest common divisor of a and b is obtained with the command `gcd(a, b)`, where in our first uses, a and b are integers. Later, a and b can be other objects with a notion of divisibility and "greatness," such as polynomials. For example,

```
gcd(2776, 2452)
```

4

We can use the `gcd` command to determine if a pair of integers are relatively prime.

```
a = 31049
b = 2105
gcd(a, b) == 1
```

True

```
a = 3563
b = 2947
gcd(a, b) == 1
```

False

The command `xgcd(a,b)` ("eXtended GCD") returns a triple where the first element is the greatest common divisor of a and b (as with the `gcd(a,b)` command above), but the next two elements are values of r and s such that $ra + sb = \gcd(a,b)$.

```
xgcd(633,331)
```

(1, -137, 262)

Portions of the triple can be extracted using `[ ]` ("indexing") to access the entries of the triple, starting with the first as number `0`. For example, the following should always return the result `True`, even if you change the values of `a` and `b`. Try changing the values of `a` and `b` below, to see that the result is always `True`.

```
a = 633
b = 331
extended = xgcd(a, b)
g = extended[0]
r = extended[1]
s = extended[2]
g == r*a + s*b
```

True

Studying this block of code will go a long way towards helping you get the most out of Sage's output. Note that `=` is how a value is *assigned* to a variable, while as in the last line, `==` is how we compare two items for *equality*.

Primes and Factoring

The method `.is_prime()` will determine if an integer is prime or not.

```
a = 117371
a.is_prime()
```

True

```
b = 14547073
b.is_prime()
```

False

```
b == 1597 * 9109
```

True

The command `random_prime(a, proof=True)` will generate a random prime number be-
tween 2 and a. Experiment by executing the following two compute cells several times.
(Replacing `proof=True` by `proof=False` will speed up the search, but there will be a very,
very, very small probability the result will not be prime.)

```
a = random_prime (10^21, proof = True)
a
```

424729101793542195193

```
a.is_prime ()
```

True

The command `prime_range(a, b)` returns an ordered list of all the primes from a to $b - 1$,
inclusive. For example,

```
prime_range (500, 550)
```

[503, 509, 521, 523, 541, 547]

The commands `next_prime(a)` and `previous_prime(a)` are other ways to get a single prime
number of a desired size. Give them a try below if you have an empty compute cell there
(as you will if you are reading in the Sage Notebook, or are reading the online version).
(The hash symbol, #, is used to indicate a "comment" line, which will not be evaluated by
Sage. So erase this line, or start on the one below it.)In addition to checking if integers are
prime or not, or generating prime numbers, Sage can also decompose any integer into its
prime factors, as described by the Fundamental Theorem of Arithmetic (Theorem 2.15).

```
a = 2600
a.factor ()
```

2^3 * 5^2 * 13

So $2600 = 2^3 \times 5^2 \times 13$ and this is the unique way to write 2600 as a product of prime
numbers (other than rearranging the order of the primes themselves in the product).

While Sage will print a factorization nicely, it is carried internally as a list of pairs of
integers, with each pair being a base (a prime number) and an exponent (a positive integer).
Study the following carefully, as it is another good exercise in working with Sage output in
the form of lists.

```
a = 2600
factored = a.factor ()
first_term = factored[0]
first_term
```

(2, 3)

```
second_term = factored[1]
second_term
```

(5, 2)

```
third_term = factored[2]
third_term
```

(13, 1)

```
first_prime = first_term[0]
first_prime
```

2

```
first_exponent = first_term[1]
first_exponent
```

3

The next compute cell reveals the internal version of the factorization by asking for the actual list. And we show how you could determine exactly how many terms the factorization has by using the length command, `len()`.

```
list(factored)
```

```
[(2, 3), (5, 2), (13, 1)]
```

```
len(factored)
```

3

Can you extract the next two primes, and their exponents, from `a`?

2.6 Sage Exercises

These exercises are about investigating basic properties of the integers, something we will frequently do when investigating groups. Sage worksheets have extensive capabilities for making new cells with carefully formatted text, include support for LATEX syntax to express mathematics. So when a question asks for explanation or commentary, make a new cell and communicate clearly with your audience.**1**.   Use the `next_prime()` command to construct two different 8-digit prime numbers and save them in variables named `a` and `b`.

2.   Use the `.is_prime()` method to verify that your primes `a` and `b` are really prime.

3.   Verify that 1 is the greatest common divisor of your two primes from the previous exercises.

4.   Find two integers that make a "linear combination" of your two primes equal to 1. Include a verification of your result.

5.   Determine a factorization into powers of primes for $c = 4\,598\,037\,234$.

6.   Write a compute cell that defines the same value of `c` again, and then defines a candidate divisor of `c` named `d`. The third line of the cell should return `True` if and only if `d` is a divisor of `c`. Illustrate the use of your cell by testing your code with $d = 7$ and in a new copy of the cell, testing your code with $d = 11$.

3

Groups

We begin our study of algebraic structures by investigating sets associated with single operations that satisfy certain reasonable axioms; that is, we want to define an operation on a set in a way that will generalize such familiar structures as the integers $\mathbb{Z}$ together with the single operation of addition, or invertible 2×2 matrices together with the single operation of matrix multiplication. The integers and the 2×2 matrices, together with their respective single operations, are examples of algebraic structures known as groups.

The theory of groups occupies a central position in mathematics. Modern group theory arose from an attempt to find the roots of a polynomial in terms of its coefficients. Groups now play a central role in such areas as coding theory, counting, and the study of symmetries; many areas of biology, chemistry, and physics have benefited from group theory.

3.1 Integer Equivalence Classes and Symmetries

Let us now investigate some mathematical structures that can be viewed as sets with single operations.

The Integers mod n

The integers mod n have become indispensable in the theory and applications of algebra. In mathematics they are used in cryptography, coding theory, and the detection of errors in identification codes.

We have already seen that two integers a and b are equivalent mod n if n divides $a - b$. The integers mod n also partition $\mathbb{Z}$ into n different equivalence classes; we will denote the set of these equivalence classes by $\mathbb{Z}_n$. Consider the integers modulo 12 and the corresponding partition of the integers:

$$[0] = \{\ldots, -12, 0, 12, 24, \ldots\},$$
$$[1] = \{\ldots, -11, 1, 13, 25, \ldots\},$$
$$\vdots$$
$$[11] = \{\ldots, -1, 11, 23, 35, \ldots\}.$$

When no confusion can arise, we will use $0, 1, \ldots, 11$ to indicate the equivalence classes $[0], [1], \ldots, [11]$ respectively. We can do arithmetic on $\mathbb{Z}_n$. For two integers a and b, define addition modulo n to be $(a + b) \pmod{n}$; that is, the remainder when $a + b$ is divided by n. Similarly, multiplication modulo n is defined as $(ab) \pmod{n}$, the remainder when ab is divided by n.

Example 3.1. The following examples illustrate integer arithmetic modulo n:

$$7 + 4 \equiv 1 \pmod 5 \qquad\qquad 7 \cdot 3 \equiv 1 \pmod 5$$
$$3 + 5 \equiv 0 \pmod 8 \qquad\qquad 3 \cdot 5 \equiv 7 \pmod 8$$
$$3 + 4 \equiv 7 \pmod{12} \qquad\qquad 3 \cdot 4 \equiv 0 \pmod{12}.$$

In particular, notice that it is possible that the product of two nonzero numbers modulo n can be equivalent to 0 modulo n.

Example 3.2. Most, but not all, of the usual laws of arithmetic hold for addition and multiplication in $\mathbb{Z}_n$. For instance, it is not necessarily true that there is a multiplicative inverse. Consider the multiplication table for $\mathbb{Z}_8$ in Table 3.3. Notice that 2, 4, and 6 do not have multiplicative inverses; that is, for $n = 2$, 4, or 6, there is no integer k such that $kn \equiv 1 \pmod 8$.

$\cdot$	0	1	2	3	4	5	6	7
0	0	0	0	0	0	0	0	0
1	0	1	2	3	4	5	6	7
2	0	2	4	6	0	2	4	6
3	0	3	6	1	4	7	2	5
4	0	4	0	4	0	4	0	4
5	0	5	2	7	4	1	6	3
6	0	6	4	2	0	6	4	2
7	0	7	6	5	4	3	2	1

Table 3.3: Multiplication table for $\mathbb{Z}_8$

Proposition 3.4. *Let $\mathbb{Z}_n$ be the set of equivalence classes of the integers mod n and $a, b, c \in \mathbb{Z}_n$.*

1. *Addition and multiplication are commutative:*

$$a + b \equiv b + a \pmod n$$
$$ab \equiv ba \pmod n.$$

2. *Addition and multiplication are associative:*

$$(a + b) + c \equiv a + (b + c) \pmod n$$
$$(ab)c \equiv a(bc) \pmod n.$$

3. *There are both additive and multiplicative identities:*

$$a + 0 \equiv a \pmod n$$
$$a \cdot 1 \equiv a \pmod n.$$

4. *Multiplication distributes over addition:*

$$a(b + c) \equiv ab + ac \pmod n.$$

5. *For every integer a there is an additive inverse −a:*

$$a + (-a) \equiv 0 \pmod{n}.$$

6. *Let a be a nonzero integer. Then* $\gcd(a, n) = 1$ *if and only if there exists a multiplicative inverse b for a* (mod n); *that is, a nonzero integer b such that*

$$ab \equiv 1 \pmod{n}.$$

PROOF.   We will prove (1) and (6) and leave the remaining properties to be proven in the exercises.

(1) Addition and multiplication are commutative modulo n since the remainder of $a + b$ divided by n is the same as the remainder of $b + a$ divided by n.

(6) Suppose that $\gcd(a, n) = 1$. Then there exist integers r and s such that $ar + ns = 1$. Since $ns = 1 - ar$, it must be the case that $ar \equiv 1 \pmod{n}$. Letting b be the equivalence class of r, $ab \equiv 1 \pmod{n}$.

Conversely, suppose that there exists an integer b such that $ab \equiv 1 \pmod{n}$. Then n divides $ab - 1$, so there is an integer k such that $ab - nk = 1$. Let $d = \gcd(a, n)$. Since d divides $ab - nk$, d must also divide 1; hence, $d = 1$.   □

Symmetries
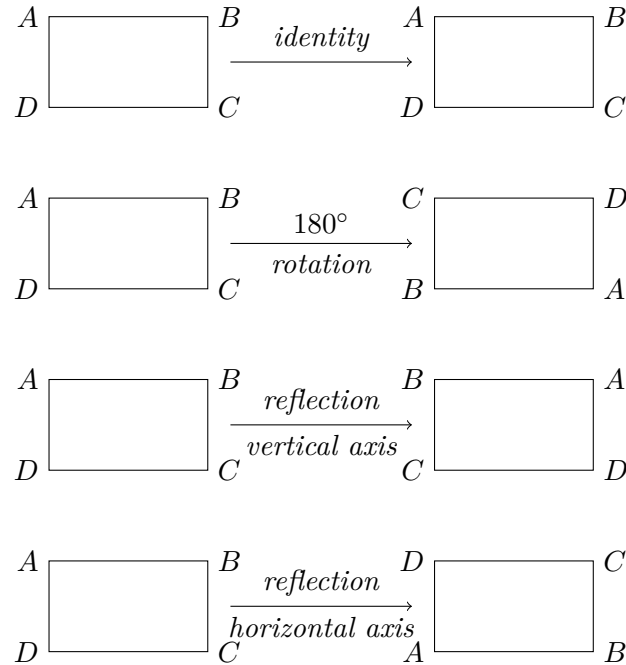
Figure 3.5: Rigid motions of a rectangle

A ***symmetry*** of a geometric figure is a rearrangement of the figure preserving the arrangement of its sides and vertices as well as its distances and angles. A map from the plane to itself preserving the symmetry of an object is called a ***rigid motion***. For example, if we look at the rectangle in Figure 3.5, it is easy to see that a rotation of 180° or 360° returns a rectangle in the plane with the same orientation as the original rectangle and the same relationship among the vertices. A reflection of the rectangle across either the vertical

axis or the horizontal axis can also be seen to be a symmetry. However, a 90° rotation in either direction cannot be a symmetry unless the rectangle is a square.

$$id = \begin{pmatrix} A & B & C \\ A & B & C \end{pmatrix}$$

$$\rho_1 = \begin{pmatrix} A & B & C \\ B & C & A \end{pmatrix}$$

$$\rho_2 = \begin{pmatrix} A & B & C \\ C & A & B \end{pmatrix}$$

$$\mu_1 = \begin{pmatrix} A & B & C \\ A & C & B \end{pmatrix}$$

$$\mu_2 = \begin{pmatrix} A & B & C \\ C & B & A \end{pmatrix}$$

$$\mu_3 = \begin{pmatrix} A & B & C \\ B & A & C \end{pmatrix}$$
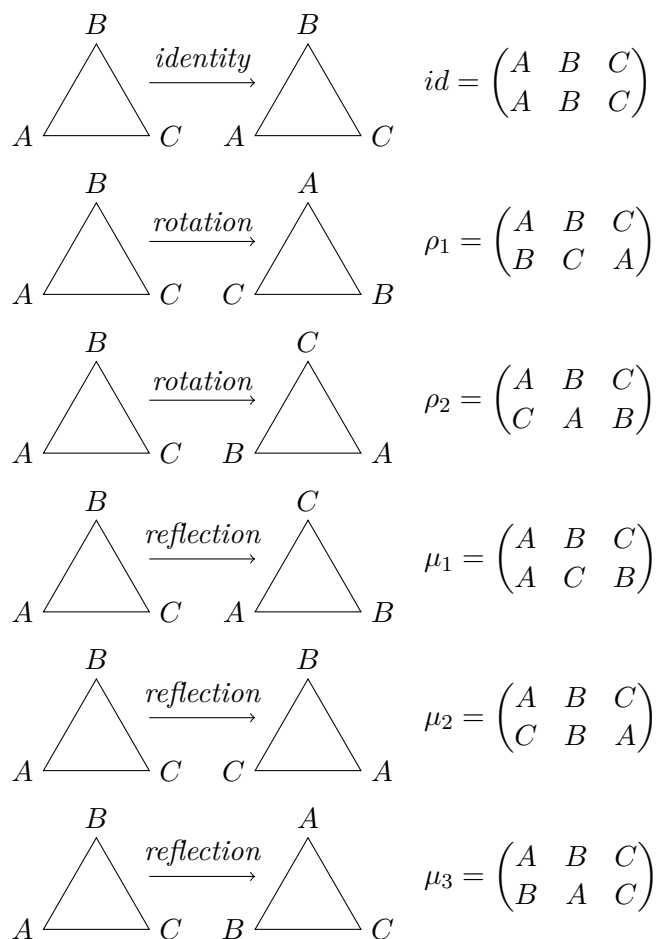
Figure 3.6: Symmetries of a triangle

Let us find the symmetries of the equilateral triangle $\triangle ABC$. To find a symmetry of $\triangle ABC$, we must first examine the permutations of the vertices A, B, and C and then ask if a permutation extends to a symmetry of the triangle. Recall that a **permutation** of a set S is a one-to-one and onto map $\pi : S \to S$. The three vertices have $3! = 6$ permutations, so the triangle has at most six symmetries. To see that there are six permutations, observe there are three different possibilities for the first vertex, and two for the second, and the remaining vertex is determined by the placement of the first two. So we have $3 \cdot 2 \cdot 1 = 3! = 6$ different arrangements. To denote the permutation of the vertices of an equilateral triangle that sends A to B, B to C, and C to A, we write the array

$$\begin{pmatrix} A & B & C \\ B & C & A \end{pmatrix}.$$

Notice that this particular permutation corresponds to the rigid motion of rotating the triangle by 120° in a clockwise direction. In fact, every permutation gives rise to a symmetry of the triangle. All of these symmetries are shown in Figure 3.6.

A natural question to ask is what happens if one motion of the triangle $\triangle ABC$ is followed by another. Which symmetry is $\mu_1 \rho_1$; that is, what happens when we do the permutation ρ_1 and then the permutation μ_1? *Remember that we are composing functions*

here. Although we usually multiply left to right, we compose functions right to left. We have

$$(\mu_1\rho_1)(A) = \mu_1(\rho_1(A)) = \mu_1(B) = C$$
$$(\mu_1\rho_1)(B) = \mu_1(\rho_1(B)) = \mu_1(C) = B$$
$$(\mu_1\rho_1)(C) = \mu_1(\rho_1(C)) = \mu_1(A) = A.$$

This is the same symmetry as μ_2. Suppose we do these motions in the opposite order, ρ_1 then μ_1. It is easy to determine that this is the same as the symmetry μ_3; hence, $\rho_1\mu_1 \neq \mu_1\rho_1$. A multiplication table for the symmetries of an equilateral triangle $\triangle ABC$ is given in Table 3.7.

Notice that in the multiplication table for the symmetries of an equilateral triangle, for every motion of the triangle α there is another motion β such that $\alpha\beta = \mathrm{id}$; that is, for every motion there is another motion that takes the triangle back to its original orientation.

$\circ$	id	ρ_1	ρ_2	μ_1	μ_2	μ_3
id	id	ρ_1	ρ_2	μ_1	μ_2	μ_3
ρ_1	ρ_1	ρ_2	id	μ_3	μ_1	μ_2
ρ_2	ρ_2	id	ρ_1	μ_2	μ_3	μ_1
μ_1	μ_1	μ_2	μ_3	id	ρ_1	ρ_2
μ_2	μ_2	μ_3	μ_1	ρ_2	id	ρ_1
μ_3	μ_3	μ_1	μ_2	ρ_1	ρ_2	id

Table 3.7: Symmetries of an equilateral triangle

3.2   Definitions and Examples

The integers mod n and the symmetries of a triangle or a rectangle are examples of groups. A **binary operation** or **law of composition** on a set G is a function $G \times G \to G$ that assigns to each pair $(a, b) \in G \times G$ a unique element $a \circ b$, or ab in G, called the composition of a and b. A **group** $(G, \circ)$ is a set G together with a law of composition $(a, b) \mapsto a \circ b$ that satisfies the following axioms.

- The law of composition is **associative**. That is,

$$(a \circ b) \circ c = a \circ (b \circ c)$$

  for $a, b, c \in G$.

- There exists an element $e \in G$, called the **identity element**, such that for any element $a \in G$

$$e \circ a = a \circ e = a.$$

- For each element $a \in G$, there exists an **inverse element** in G, denoted by a^{-1}, such that

$$a \circ a^{-1} = a^{-1} \circ a = e.$$

A group G with the property that $a \circ b = b \circ a$ for all $a, b \in G$ is called **abelian** or **commutative**. Groups not satisfying this property are said to be **nonabelian** or **noncommutative**.

Example 3.8. The integers $\mathbb{Z} = \{\ldots, -1, 0, 1, 2, \ldots\}$ form a group under the operation of addition. The binary operation on two integers $m, n \in \mathbb{Z}$ is just their sum. Since the integers under addition already have a well-established notation, we will use the operator $+$ instead of $\circ$; that is, we shall write $m + n$ instead of $m \circ n$. The identity is 0, and the inverse of $n \in \mathbb{Z}$ is written as $-n$ instead of n^{-1}. Notice that the set of integers under addition have the additional property that $m + n = n + m$ and therefore form an abelian group.

Most of the time we will write ab instead of $a \circ b$; however, if the group already has a natural operation such as addition in the integers, we will use that operation. That is, if we are adding two integers, we still write $m + n$, $-n$ for the inverse, and 0 for the identity as usual. We also write $m - n$ instead of $m + (-n)$.

It is often convenient to describe a group in terms of an addition or multiplication table. Such a table is called a **Cayley table**.

Example 3.9. The integers mod n form a group under addition modulo n. Consider $\mathbb{Z}_5$, consisting of the equivalence classes of the integers 0, 1, 2, 3, and 4. We define the group operation on $\mathbb{Z}_5$ by modular addition. We write the binary operation on the group additively; that is, we write $m + n$. The element 0 is the identity of the group and each element in $\mathbb{Z}_5$ has an inverse. For instance, $2 + 3 = 3 + 2 = 0$. Table 3.10 is a Cayley table for $\mathbb{Z}_5$. By Proposition 3.4, $\mathbb{Z}_n = \{0, 1, \ldots, n-1\}$ is a group under the binary operation of addition mod n.

$+$	0	1	2	3	4
0	0	1	2	3	4
1	1	2	3	4	0
2	2	3	4	0	1
3	3	4	0	1	2
4	4	0	1	2	3

Table 3.10: Cayley table for $(\mathbb{Z}_5, +)$

Example 3.11. Not every set with a binary operation is a group. For example, if we let modular multiplication be the binary operation on $\mathbb{Z}_n$, then $\mathbb{Z}_n$ fails to be a group. The element 1 acts as a group identity since $1 \cdot k = k \cdot 1 = k$ for any $k \in \mathbb{Z}_n$; however, a multiplicative inverse for 0 does not exist since $0 \cdot k = k \cdot 0 = 0$ for every k in $\mathbb{Z}_n$. Even if we consider the set $\mathbb{Z}_n \setminus \{0\}$, we still may not have a group. For instance, let $2 \in \mathbb{Z}_6$. Then 2 has no multiplicative inverse since

$$0 \cdot 2 = 0 \qquad 1 \cdot 2 = 2$$
$$2 \cdot 2 = 4 \qquad 3 \cdot 2 = 0$$
$$4 \cdot 2 = 2 \qquad 5 \cdot 2 = 4.$$

By Proposition 3.4, every nonzero k does have an inverse in $\mathbb{Z}_n$ if k is relatively prime to n. Denote the set of all such nonzero elements in $\mathbb{Z}_n$ by $U(n)$. Then $U(n)$ is a group called the **group of units** of $\mathbb{Z}_n$. Table 3.12 is a Cayley table for the group $U(8)$.

·	1	3	5	7
1	1	3	5	7
3	3	1	7	5
5	5	7	1	3
7	7	5	3	1

Table 3.12: Multiplication table for $U(8)$

Example 3.13. The symmetries of an equilateral triangle described in Section 3.1 form a nonabelian group. As we observed, it is not necessarily true that $\alpha\beta = \beta\alpha$ for two symmetries α and β. Using Table 3.7, which is a Cayley table for this group, we can easily check that the symmetries of an equilateral triangle are indeed a group. We will denote this group by either S_3 or D_3, for reasons that will be explained later.

Example 3.14. We use $\mathbb{M}_2(\mathbb{R})$ to denote the set of all 2×2 matrices. Let $GL_2(\mathbb{R})$ be the subset of $\mathbb{M}_2(\mathbb{R})$ consisting of invertible matrices; that is, a matrix

$$A = \begin{pmatrix} a & b \\ c & d \end{pmatrix}$$

is in $GL_2(\mathbb{R})$ if there exists a matrix A^{-1} such that $AA^{-1} = A^{-1}A = I$, where I is the 2×2 identity matrix. For A to have an inverse is equivalent to requiring that the determinant of A be nonzero; that is, $\det A = ad - bc \neq 0$. The set of invertible matrices forms a group called the **general linear group**. The identity of the group is the identity matrix

$$I = \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}.$$

The inverse of $A \in GL_2(\mathbb{R})$ is

$$A^{-1} = \frac{1}{ad - bc} \begin{pmatrix} d & -b \\ -c & a \end{pmatrix}.$$

The product of two invertible matrices is again invertible. Matrix multiplication is associative, satisfying the other group axiom. For matrices it is not true in general that $AB = BA$; hence, $GL_2(\mathbb{R})$ is another example of a nonabelian group.

Example 3.15. Let

$$1 = \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix} \qquad I = \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}$$

$$J = \begin{pmatrix} 0 & i \\ i & 0 \end{pmatrix} \qquad K = \begin{pmatrix} i & 0 \\ 0 & -i \end{pmatrix},$$

where $i^2 = -1$. Then the relations $I^2 = J^2 = K^2 = -1$, $IJ = K$, $JK = I$, $KI = J$, $JI = -K$, $KJ = -I$, and $IK = -J$ hold. The set $Q_8 = \{\pm 1, \pm I, \pm J, \pm K\}$ is a group called the **quaternion group**. Notice that Q_8 is noncommutative.

Example 3.16. Let $\mathbb{C}^*$ be the set of nonzero complex numbers. Under the operation of multiplication $\mathbb{C}^*$ forms a group. The identity is 1. If $z = a + bi$ is a nonzero complex number, then

$$z^{-1} = \frac{a - bi}{a^2 + b^2}$$

is the inverse of z. It is easy to see that the remaining group axioms hold.

A group is ***finite***, or has ***finite order***, if it contains a finite number of elements; otherwise, the group is said to be ***infinite*** or to have ***infinite order***. The ***order*** of a finite group is the number of elements that it contains. If G is a group containing n elements, we write $|G| = n$. The group $\mathbb{Z}_5$ is a finite group of order 5; the integers $\mathbb{Z}$ form an infinite group under addition, and we sometimes write $|\mathbb{Z}| = \infty$.

Basic Properties of Groups

Proposition 3.17. *The identity element in a group G is unique; that is, there exists only one element $e \in G$ such that $eg = ge = g$ for all $g \in G$.*

PROOF. Suppose that e and e' are both identities in G. Then $eg = ge = g$ and $e'g = ge' = g$ for all $g \in G$. We need to show that $e = e'$. If we think of e as the identity, then $ee' = e'$; but if e' is the identity, then $ee' = e$. Combining these two equations, we have $e = ee' = e'$. □

Inverses in a group are also unique. If g' and g'' are both inverses of an element g in a group G, then $gg' = g'g = e$ and $gg'' = g''g = e$. We want to show that $g' = g''$, but $g' = g'e = g'(gg'') = (g'g)g'' = eg'' = g''$. We summarize this fact in the following proposition.

Proposition 3.18. *If g is any element in a group G, then the inverse of g, denoted by g^{-1}, is unique.*

Proposition 3.19. *Let G be a group. If $a, b \in G$, then $(ab)^{-1} = b^{-1}a^{-1}$.*

PROOF. Let $a, b \in G$. Then $abb^{-1}a^{-1} = aea^{-1} = aa^{-1} = e$. Similarly, $b^{-1}a^{-1}ab = e$. But by the previous proposition, inverses are unique; hence, $(ab)^{-1} = b^{-1}a^{-1}$. □

Proposition 3.20. *Let G be a group. For any $a \in G$, $(a^{-1})^{-1} = a$.*

PROOF. Observe that $a^{-1}(a^{-1})^{-1} = e$. Consequently, multiplying both sides of this equation by a, we have

$$(a^{-1})^{-1} = e(a^{-1})^{-1} = aa^{-1}(a^{-1})^{-1} = ae = a.$$ □

It makes sense to write equations with group elements and group operations. If a and b are two elements in a group G, does there exist an element $x \in G$ such that $ax = b$? If such an x does exist, is it unique? The following proposition answers both of these questions positively.

Proposition 3.21. *Let G be a group and a and b be any two elements in G. Then the equations $ax = b$ and $xa = b$ have unique solutions in G.*

PROOF. Suppose that $ax = b$. We must show that such an x exists. We can multiply both sides of $ax = b$ by a^{-1} to find $x = ex = a^{-1}ax = a^{-1}b$.

To show uniqueness, suppose that x_1 and x_2 are both solutions of $ax = b$; then $ax_1 = b = ax_2$. So $x_1 = a^{-1}ax_1 = a^{-1}ax_2 = x_2$. The proof for the existence and uniqueness of the solution of $xa = b$ is similar. □

Proposition 3.22. *If G is a group and $a, b, c \in G$, then $ba = ca$ implies $b = c$ and $ab = ac$ implies $b = c$.*

This proposition tells us that the ***right and left cancellation laws*** are true in groups. We leave the proof as an exercise.

We can use exponential notation for groups just as we do in ordinary algebra. If G is a group and $g \in G$, then we define $g^0 = e$. For $n \in \mathbb{N}$, we define

$$g^n = \underbrace{g \cdot g \cdots g}_{n \text{ times}}$$

and

$$g^{-n} = \underbrace{g^{-1} \cdot g^{-1} \cdots g^{-1}}_{n \text{ times}}.$$

Theorem 3.23. *In a group, the usual laws of exponents hold; that is, for all $g, h \in G$,*

1. *$g^m g^n = g^{m+n}$ for all $m, n \in \mathbb{Z}$;*

2. *$(g^m)^n = g^{mn}$ for all $m, n \in \mathbb{Z}$;*

3. *$(gh)^n = (h^{-1}g^{-1})^{-n}$ for all $n \in \mathbb{Z}$. Furthermore, if G is abelian, then $(gh)^n = g^n h^n$.*

We will leave the proof of this theorem as an exercise. Notice that $(gh)^n \neq g^n h^n$ in general, since the group may not be abelian. If the group is $\mathbb{Z}$ or $\mathbb{Z}_n$, we write the group operation additively and the exponential operation multiplicatively; that is, we write ng instead of g^n. The laws of exponents now become

1. $mg + ng = (m+n)g$ for all $m, n \in \mathbb{Z}$;

2. $m(ng) = (mn)g$ for all $m, n \in \mathbb{Z}$;

3. $m(g + h) = mg + mh$ for all $n \in \mathbb{Z}$.

It is important to realize that the last statement can be made only because $\mathbb{Z}$ and $\mathbb{Z}_n$ are commutative groups.

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Historical Note

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Although the first clear axiomatic definition of a group was not given until the late 1800s, group-theoretic methods had been employed before this time in the development of many areas of mathematics, including geometry and the theory of algebraic equations.

Joseph-Louis Lagrange used group-theoretic methods in a 1770–1771 memoir to study methods of solving polynomial equations. Later, Évariste Galois (1811–1832) succeeded in developing the mathematics necessary to determine exactly which polynomial equations could be solved in terms of the polynomials'coefficients. Galois' primary tool was group theory.

The study of geometry was revolutionized in 1872 when Felix Klein proposed that geometric spaces should be studied by examining those properties that are invariant under a transformation of the space. Sophus Lie, a contemporary of Klein, used group theory to study solutions of partial differential equations. One of the first modern treatments of group theory appeared in William Burnside's *The Theory of Groups of Finite Order* [1], first published in 1897.

3.3 Subgroups

Definitions and Examples

Sometimes we wish to investigate smaller groups sitting inside a larger group. The set of even integers $2\mathbb{Z} = \{\ldots, -2, 0, 2, 4, \ldots\}$ is a group under the operation of addition. This smaller group sits naturally inside of the group of integers under addition. We define a **subgroup** H of a group G to be a subset H of G such that when the group operation of G is restricted to H, H is a group in its own right. Observe that every group G with at least two elements will always have at least two subgroups, the subgroup consisting of the identity element alone and the entire group itself. The subgroup $H = \{e\}$ of a group G is called the **trivial subgroup**. A subgroup that is a proper subset of G is called a **proper subgroup**. In many of the examples that we have investigated up to this point, there exist other subgroups besides the trivial and improper subgroups.

Example 3.24. Consider the set of nonzero real numbers, $\mathbb{R}^*$, with the group operation of multiplication. The identity of this group is 1 and the inverse of any element $a \in \mathbb{R}^*$ is just $1/a$. We will show that

$$\mathbb{Q}^* = \{p/q : p \text{ and } q \text{ are nonzero integers}\}$$

is a subgroup of $\mathbb{R}^*$. The identity of $\mathbb{R}^*$ is 1; however, $1 = 1/1$ is the quotient of two nonzero integers. Hence, the identity of $\mathbb{R}^*$ is in $\mathbb{Q}^*$. Given two elements in $\mathbb{Q}^*$, say p/q and r/s, their product pr/qs is also in $\mathbb{Q}^*$. The inverse of any element $p/q \in \mathbb{Q}^*$ is again in $\mathbb{Q}^*$ since $(p/q)^{-1} = q/p$. Since multiplication in $\mathbb{R}^*$ is associative, multiplication in $\mathbb{Q}^*$ is associative.

Example 3.25. Recall that $\mathbb{C}^*$ is the multiplicative group of nonzero complex numbers. Let $H = \{1, -1, i, -i\}$. Then H is a subgroup of $\mathbb{C}^*$. It is quite easy to verify that H is a group under multiplication and that $H \subset \mathbb{C}^*$.

Example 3.26. Let $SL_2(\mathbb{R})$ be the subset of $GL_2(\mathbb{R})$ consisting of matrices of determinant one; that is, a matrix

$$A = \begin{pmatrix} a & b \\ c & d \end{pmatrix}$$

is in $SL_2(\mathbb{R})$ exactly when $ad - bc = 1$. To show that $SL_2(\mathbb{R})$ is a subgroup of the general linear group, we must show that it is a group under matrix multiplication. The 2×2 identity matrix is in $SL_2(\mathbb{R})$, as is the inverse of the matrix A:

$$A^{-1} = \begin{pmatrix} d & -b \\ -c & a \end{pmatrix}.$$

It remains to show that multiplication is closed; that is, that the product of two matrices of determinant one also has determinant one. We will leave this task as an exercise. The group $SL_2(\mathbb{R})$ is called the **special linear group**.

Example 3.27. It is important to realize that a subset H of a group G can be a group without being a subgroup of G. For H to be a subgroup of G it must inherit G's binary operation. The set of all 2×2 matrices, $\mathbb{M}_2(\mathbb{R})$, forms a group under the operation of addition. The 2×2 general linear group is a subset of $\mathbb{M}_2(\mathbb{R})$ and is a group under matrix multiplication, but it is not a subgroup of $\mathbb{M}_2(\mathbb{R})$. If we add two invertible matrices, we do not necessarily obtain another invertible matrix. Observe that

$$\begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix} + \begin{pmatrix} -1 & 0 \\ 0 & -1 \end{pmatrix} = \begin{pmatrix} 0 & 0 \\ 0 & 0 \end{pmatrix},$$

but the zero matrix is not in $GL_2(\mathbb{R})$.

Example 3.28. One way of telling whether or not two groups are the same is by examining their subgroups. Other than the trivial subgroup and the group itself, the group $\mathbb{Z}_4$ has a single subgroup consisting of the elements 0 and 2. From the group $\mathbb{Z}_2$, we can form another group of four elements as follows. As a set this group is $\mathbb{Z}_2 \times \mathbb{Z}_2$. We perform the group operation coordinatewise; that is, $(a, b) + (c, d) = (a + c, b + d)$. Table 3.29 is an addition table for $\mathbb{Z}_2 \times \mathbb{Z}_2$. Since there are three nontrivial proper subgroups of $\mathbb{Z}_2 \times \mathbb{Z}_2$, $H_1 = \{(0,0), (0,1)\}$, $H_2 = \{(0,0), (1,0)\}$, and $H_3 = \{(0,0), (1,1)\}$, $\mathbb{Z}_4$ and $\mathbb{Z}_2 \times \mathbb{Z}_2$ must be different groups.

$+$	$(0,0)$	$(0,1)$	$(1,0)$	$(1,1)$
$(0,0)$	$(0,0)$	$(0,1)$	$(1,0)$	$(1,1)$
$(0,1)$	$(0,1)$	$(0,0)$	$(1,1)$	$(1,0)$
$(1,0)$	$(1,0)$	$(1,1)$	$(0,0)$	$(0,1)$
$(1,1)$	$(1,1)$	$(1,0)$	$(0,1)$	$(0,0)$

Table 3.29: Addition table for $\mathbb{Z}_2 \times \mathbb{Z}_2$

Some Subgroup Theorems

Let us examine some criteria for determining exactly when a subset of a group is a subgroup.

Proposition 3.30. *A subset H of G is a subgroup if and only if it satisfies the following conditions.*

1. *The identity e of G is in H.*

2. *If $h_1, h_2 \in H$, then $h_1 h_2 \in H$.*

3. *If $h \in H$, then $h^{-1} \in H$.*

PROOF.   First suppose that H is a subgroup of G. We must show that the three conditions hold. Since H is a group, it must have an identity e_H. We must show that $e_H = e$, where e is the identity of G. We know that $e_H e_H = e_H$ and that $e e_H = e_H e = e_H$; hence, $e e_H = e_H e_H$. By right-hand cancellation, $e = e_H$. The second condition holds since a subgroup H is a group. To prove the third condition, let $h \in H$. Since H is a group, there is an element $h' \in H$ such that $h h' = h' h = e$. By the uniqueness of the inverse in G, $h' = h^{-1}$.

Conversely, if the three conditions hold, we must show that H is a group under the same operation as G; however, these conditions plus the associativity of the binary operation are exactly the axioms stated in the definition of a group. $\square$

Proposition 3.31. *Let H be a subset of a group G. Then H is a subgroup of G if and only if $H \neq \emptyset$, and whenever $g, h \in H$ then gh^{-1} is in H.*

PROOF.   First assume that H is a subgroup of G. We wish to show that $gh^{-1} \in H$ whenever g and h are in H. Since h is in H, its inverse h^{-1} must also be in H. Because of the closure of the group operation, $gh^{-1} \in H$.

Conversely, suppose that $H \subset G$ such that $H \neq \emptyset$ and $gh^{-1} \in H$ whenever $g, h \in H$. If $g \in H$, then $gg^{-1} = e$ is in H. If $g \in H$, then $eg^{-1} = g^{-1}$ is also in H. Now let $h_1, h_2 \in H$. We must show that their product is also in H. However, $h_1(h_2^{-1})^{-1} = h_1 h_2 \in H$. Hence, H is a subgroup of G. $\square$

3.4 Exercises

1.   Find all $x \in \mathbb{Z}$ satisfying each of the following equations.

(a) $3x \equiv 2 \pmod{7}$ (d) $9x \equiv 3 \pmod{5}$

(b) $5x + 1 \equiv 13 \pmod{23}$ (e) $5x \equiv 1 \pmod{6}$

(c) $5x + 1 \equiv 13 \pmod{26}$ (f) $3x \equiv 1 \pmod{6}$

2.   Which of the following multiplication tables defined on the set $G = \{a, b, c, d\}$ form a group? Support your answer in each case.

(a)

$\circ$	a	b	c	d
a	a	c	d	a
b	b	b	c	d
c	c	d	a	b
d	d	a	b	c

(c)

$\circ$	a	b	c	d
a	a	b	c	d
b	b	c	d	a
c	c	d	a	b
d	d	a	b	c

(b)

$\circ$	a	b	c	d
a	a	b	c	d
b	b	a	d	c
c	c	d	a	b
d	d	c	b	a

(d)

$\circ$	a	b	c	d
a	a	b	c	d
b	b	a	c	d
c	c	b	a	d
d	d	d	b	c

3.   Write out Cayley tables for groups formed by the symmetries of a rectangle and for $(\mathbb{Z}_4, +)$. How many elements are in each group? Are the groups the same? Why or why not?

4.   Describe the symmetries of a rhombus and prove that the set of symmetries forms a group. Give Cayley tables for both the symmetries of a rectangle and the symmetries of a rhombus. Are the symmetries of a rectangle and those of a rhombus the same?

5.   Describe the symmetries of a square and prove that the set of symmetries is a group. Give a Cayley table for the symmetries. How many ways can the vertices of a square be permuted? Is each permutation necessarily a symmetry of the square? The symmetry group of the square is denoted by D_4.

6.   Give a multiplication table for the group $U(12)$.

7.   Let $S = \mathbb{R} \setminus \{-1\}$ and define a binary operation on S by $a * b = a + b + ab$. Prove that $(S, *)$ is an abelian group.

8.   Give an example of two elements A and B in $GL_2(\mathbb{R})$ with $AB \neq BA$.

9.   Prove that the product of two matrices in $SL_2(\mathbb{R})$ has determinant one.

10.   Prove that the set of matrices of the form

$$\begin{pmatrix} 1 & x & y \\ 0 & 1 & z \\ 0 & 0 & 1 \end{pmatrix}$$

is a group under matrix multiplication. This group, known as the **Heisenberg group**, is important in quantum physics. Matrix multiplication in the Heisenberg group is defined by

$$
\begin{pmatrix} 1 & x & y \\ 0 & 1 & z \\ 0 & 0 & 1 \end{pmatrix} \begin{pmatrix} 1 & x' & y' \\ 0 & 1 & z' \\ 0 & 0 & 1 \end{pmatrix} = \begin{pmatrix} 1 & x + x' & y + y' + xz' \\ 0 & 1 & z + z' \\ 0 & 0 & 1 \end{pmatrix}.
$$

11.   Prove that $\det(AB) = \det(A)\det(B)$ in $GL_2(\mathbb{R})$. Use this result to show that the binary operation in the group $GL_2(\mathbb{R})$ is closed; that is, if A and B are in $GL_2(\mathbb{R})$, then $AB \in GL_2(\mathbb{R})$.

12.   Let $\mathbb{Z}_2^n = \{(a_1, a_2, \ldots, a_n) : a_i \in \mathbb{Z}_2\}$. Define a binary operation on $\mathbb{Z}_2^n$ by

$$(a_1, a_2, \ldots, a_n) + (b_1, b_2, \ldots, b_n) = (a_1 + b_1, a_2 + b_2, \ldots, a_n + b_n).$$

Prove that $\mathbb{Z}_2^n$ is a group under this operation. This group is important in algebraic coding theory.

13.   Show that $\mathbb{R}^* = \mathbb{R} \setminus \{0\}$ is a group under the operation of multiplication.

14.   Given the groups $\mathbb{R}^*$ and $\mathbb{Z}$, let $G = \mathbb{R}^* \times \mathbb{Z}$. Define a binary operation $\circ$ on G by $(a, m) \circ (b, n) = (ab, m + n)$. Show that G is a group under this operation.

15.   Prove or disprove that every group containing six elements is abelian.

16.   Give a specific example of some group G and elements $g, h \in G$ where $(gh)^n \neq g^n h^n$.

17.   Give an example of three different groups with eight elements. Why are the groups different?

18.   Show that there are $n!$ permutations of a set containing n items.

19.   Show that

$$0 + a \equiv a + 0 \equiv a \pmod{n}$$

for all $a \in \mathbb{Z}_n$.

20.   Prove that there is a multiplicative identity for the integers modulo n:

$$a \cdot 1 \equiv a \pmod{n}.$$

21.   For each $a \in \mathbb{Z}_n$ find an element $b \in \mathbb{Z}_n$ such that

$$a + b \equiv b + a \equiv 0 \pmod{n}.$$

22.   Show that addition and multiplication mod n are well defined operations. That is, show that the operations do not depend on the choice of the representative from the equivalence classes mod n.

23.   Show that addition and multiplication mod n are associative operations.

24.   Show that multiplication distributes over addition modulo n:

$$a(b + c) \equiv ab + ac \pmod{n}.$$

25.   Let a and b be elements in a group G. Prove that $ab^n a^{-1} = (aba^{-1})^n$ for $n \in \mathbb{Z}$.

26.   Let $U(n)$ be the group of units in $\mathbb{Z}_n$. If $n > 2$, prove that there is an element $k \in U(n)$ such that $k^2 = 1$ and $k \neq 1$.

27.   Prove that the inverse of $g_1 g_2 \cdots g_n$ is $g_n^{-1} g_{n-1}^{-1} \cdots g_1^{-1}$.

28.   Prove the remainder of Proposition 3.21: if G is a group and $a, b \in G$, then the equation $xa = b$ has a unique solution in G.

29.   Prove Theorem 3.23.

30.   Prove the right and left cancellation laws for a group G; that is, show that in the group G, $ba = ca$ implies $b = c$ and $ab = ac$ implies $b = c$ for elements $a, b, c \in G$.

31.   Show that if $a^2 = e$ for all elements a in a group G, then G must be abelian.

32.   Show that if G is a finite group of even order, then there is an $a \in G$ such that a is not the identity and $a^2 = e$.

33.   Let G be a group and suppose that $(ab)^2 = a^2b^2$ for all a and b in G. Prove that G is an abelian group.

34.   Find all the subgroups of $\mathbb{Z}_3 \times \mathbb{Z}_3$. Use this information to show that $\mathbb{Z}_3 \times \mathbb{Z}_3$ is not the same group as $\mathbb{Z}_9$. (See Example 3.28 for a short description of the product of groups.)

35.   Find all the subgroups of the symmetry group of an equilateral triangle.

36.   Compute the subgroups of the symmetry group of a square.

37.   Let $H = \{2^k : k \in \mathbb{Z}\}$. Show that H is a subgroup of $\mathbb{Q}^*$.

38.   Let $n = 0, 1, 2, \ldots$ and $n\mathbb{Z} = \{nk : k \in \mathbb{Z}\}$. Prove that $n\mathbb{Z}$ is a subgroup of $\mathbb{Z}$. Show that these subgroups are the only subgroups of $\mathbb{Z}$.

39.   Let $\mathbb{T} = \{z \in \mathbb{C}^* : |z| = 1\}$. Prove that $\mathbb{T}$ is a subgroup of $\mathbb{C}^*$.

40.   Let G consist of the 2×2 matrices of the form

$$\begin{pmatrix} \cos\theta & -\sin\theta \\ \sin\theta & \cos\theta \end{pmatrix},$$

where $\theta \in \mathbb{R}$. Prove that G is a subgroup of $SL_2(\mathbb{R})$.

41.   Prove that

$$G = \{a + b\sqrt{2} : a, b \in \mathbb{Q} \text{ and } a \text{ and } b \text{ are not both zero}\}$$

is a subgroup of $\mathbb{R}^*$ under the group operation of multiplication.

42.   Let G be the group of 2×2 matrices under addition and

$$H = \left\{ \begin{pmatrix} a & b \\ c & d \end{pmatrix} : a + d = 0 \right\}.$$

Prove that H is a subgroup of G.

43.   Prove or disprove: $SL_2(\mathbb{Z})$, the set of 2×2 matrices with integer entries and determinant one, is a subgroup of $SL_2(\mathbb{R})$.

44.   List the subgroups of the quaternion group, Q_8.

45.   Prove that the intersection of two subgroups of a group G is also a subgroup of G.

46.   Prove or disprove: If H and K are subgroups of a group G, then $H \cup K$ is a subgroup of G.

47.   Prove or disprove: If H and K are subgroups of a group G, then $HK = \{hk : h \in H \text{ and } k \in K\}$ is a subgroup of G. What if G is abelian?

48.   Let G be a group and $g \in G$. Show that

$$Z(G) = \{x \in G : gx = xg \text{ for all } g \in G\}$$

is a subgroup of G. This subgroup is called the **center** of G.

49.   Let a and b be elements of a group G. If $a^4b = ba$ and $a^3 = e$, prove that $ab = ba$.

50.   Give an example of an infinite group in which every nontrivial subgroup is infinite.

51.   If $xy = x^{-1}y^{-1}$ for all x and y in G, prove that G must be abelian.

52.   Prove or disprove: Every proper subgroup of an nonabelian group is nonabelian.

53.   Let H be a subgroup of G and

$$C(H) = \{g \in G : gh = hg \text{ for all } h \in H\}.$$

Prove $C(H)$ is a subgroup of G. This subgroup is called the ***centralizer*** of H in G.

54.   Let H be a subgroup of G. If $g \in G$, show that $gHg^{-1} = \{ghg^{-1} : h \in H\}$ is also a subgroup of G.

3.5   Additional Exercises: Detecting Errors

1. **UPC Symbols.**   Universal Product Code (UPC) symbols are found on most products in grocery and retail stores. The UPC symbol is a 12-digit code identifying the manufacturer of a product and the product itself (Figure 3.32). The first 11 digits contain information about the product; the twelfth digit is used for error detection. If $d_1 d_2 \cdots d_{12}$ is a valid UPC number, then

$$3 \cdot d_1 + 1 \cdot d_2 + 3 \cdot d_3 + \cdots + 3 \cdot d_{11} + 1 \cdot d_{12} \equiv 0 \pmod{10}.$$

(a) Show that the UPC number 0-50000-30042-6, which appears in Figure 3.32, is a valid UPC number.

(b) Show that the number 0-50000-30043-6 is not a valid UPC number.

(c) Write a formula to calculate the check digit, d_{12}, in the UPC number.

(d) The UPC error detection scheme can detect most transposition errors; that is, it can determine if two digits have been interchanged. Show that the transposition error 0-05000-30042-6 is not detected. Find a transposition error that is detected. Can you find a general rule for the types of transposition errors that can be detected?

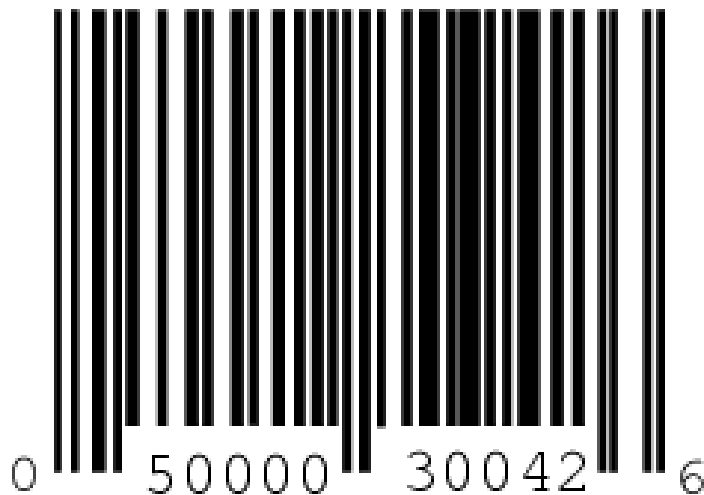(e) Write a program that will determine whether or not a UPC number is valid.

Figure 3.32: A UPC code

2.   It is often useful to use an inner product notation for this type of error detection scheme; hence, we will use the notion

$$(d_1, d_2, \ldots, d_k) \cdot (w_1, w_2, \ldots, w_k) \equiv 0 \pmod{n}$$

to mean

$$d_1 w_1 + d_2 w_2 + \cdots + d_k w_k \equiv 0 \pmod{n}.$$

Suppose that $(d_1, d_2, \ldots, d_k) \cdot (w_1, w_2, \ldots, w_k) \equiv 0 \pmod{n}$ is an error detection scheme for the k-digit identification number $d_1 d_2 \cdots d_k$, where $0 \leq d_i < n$. Prove that all single-digit errors are detected if and only if $\gcd(w_i, n) = 1$ for $1 \leq i \leq k$.

3.   Let $(d_1, d_2, \ldots, d_k) \cdot (w_1, w_2, \ldots, w_k) \equiv 0 \pmod{n}$ be an error detection scheme for the k-digit identification number $d_1 d_2 \cdots d_k$, where $0 \leq d_i < n$. Prove that all transposition errors of two digits d_i and d_j are detected if and only if $\gcd(w_i - w_j, n) = 1$ for i and j between 1 and k.

4. **ISBN Codes.**   Every book has an International Standard Book Number (ISBN) code. This is a 10-digit code indicating the book's publisher and title. The tenth digit is a check digit satisfying

$$(d_1, d_2, \ldots, d_{10}) \cdot (10, 9, \ldots, 1) \equiv 0 \pmod{11}.$$

One problem is that d_{10} might have to be a 10 to make the inner product zero; in this case, 11 digits would be needed to make this scheme work. Therefore, the character X is used for the eleventh digit. So ISBN 3-540-96035-X is a valid ISBN code.

(a) Is ISBN 0-534-91500-0 a valid ISBN code? What about ISBN 0-534-91700-0 and ISBN 0-534-19500-0?

(b) Does this method detect all single-digit errors? What about all transposition errors?

(c) How many different ISBN codes are there?

(d) Write a computer program that will calculate the check digit for the first nine digits of an ISBN code.

(e) A publisher has houses in Germany and the United States. Its German prefix is 3-540. If its United States prefix will be 0-abc, find abc such that the rest of the ISBN code will be the same for a book printed in Germany and in the United States. Under the ISBN coding method the first digit identifies the language; German is 3 and English is 0. The next group of numbers identifies the publisher, and the last group identifies the specific book.

References and Suggested Readings

[1]   Burnside, W. *Theory of Groups of Finite Order.* 2nd ed. Cambridge University Press, Cambridge, 1911; Dover, New York, 1953. A classic. Also available at books.google.com.

[2]   Gallian, J. A. and Winters, S. "Modular Arithmetic in the Marketplace," *The American Mathematical Monthly* **95** (1988): 548–51.

[3]   Gallian, J. A. *Contemporary Abstract Algebra.* 7th ed. Brooks/Cole, Belmont, CA, 2009.

[4]   Hall, M. *Theory of Groups.* 2nd ed. American Mathematical Society, Providence, 1959.

[**5**]   Kurosh, A. E. *The Theory of Groups*, vols. I and II. American Mathematical Society, Providence, 1979.

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3.6   Sage

Many of the groups discussed in this chapter are available for study in Sage. It is important to understand that sets that form algebraic objects (groups in this chapter) are called "parents" in Sage, and elements of these objects are called, well, "elements." So every element belongs to a parent (in other words, is contained in some set). We can ask about properties of parents (finite? order? abelian?), and we can ask about properties of individual elements (identity? inverse?). In the following we will show you how to create some of these common groups and begin to explore their properties with Sage.

Integers mod n

```
Z8 = Integers(8)
Z8
```

Ring of integers modulo 8

```
Z8.list()
```

[0, 1, 2, 3, 4, 5, 6, 7]

```
a = Z8.an_element(); a
```

0

```
a.parent()
```

Ring of integers modulo 8

We would like to work with elements of Z8. If you were to type a 6 into a compute cell right now, what would you mean? The integer 6, the rational number $\frac{6}{1}$, the real number 6.00000, or the complex number $6.00000 + 0.00000i$? Or perhaps you really do want the integer 6 mod 8? Sage really has no idea what you mean or want. To make this clear, you can "coerce" 6 into Z8 with the syntax Z8(6). Without this, Sage will treat a input number like 6 as an integer, the simplest possible interpretation in some sense. Study the following carefully, where we first work with "normal" integers and then with integers mod 8.

```
a = 6
a
```

6

```
a.parent()
```

Integer Ring

```
b = 7
c = a + b; c
```

13

```
d = Z8(6)
d
```

6

```
d.parent()
```

Ring of integers modulo 8

```
e = Z8(7)
f = d+e; f
```

5

```
g = Z8(85); g
```

5

```
f == g
```

True

Z8 is a bit unusual as a first example, since it has two operations defined, both addition and multiplication, with addition forming a group, and multiplication not forming a group. Still, we can work with the additive portion, here forming the Cayley table for the addition.

```
Z8.addition_table(names='elements')
```

```
+  0 1 2 3 4 5 6 7
 +----------------
0|  0 1 2 3 4 5 6 7
1|  1 2 3 4 5 6 7 0
2|  2 3 4 5 6 7 0 1
3|  3 4 5 6 7 0 1 2
4|  4 5 6 7 0 1 2 3
5|  5 6 7 0 1 2 3 4
6|  6 7 0 1 2 3 4 5
7|  7 0 1 2 3 4 5 6
```

When n is a prime number, the multipicative structure (excluding zero), will also form a group.

The integers mod n are very important, so Sage implements both addition and multiplication together. Groups of symmetries are a better example of how Sage implements groups, since there is just one operation present.

Groups of symmetries

The symmetries of some geometric shapes are already defined in Sage, albeit with different names. They are implemented as "permutation groups" which we will begin to study carefully in Chapter 5.

Sage uses integers to label vertices, starting the count at 1, instead of letters. Elements by default are printed using "cycle notation" which we will see described carefully in Chapter 5. Here is an example, with both the mathematics and Sage. For the Sage part, we create the group of symmetries and then create the symmetry ρ_2 with coercion, followed

by outputting the element in cycle notation. Then we create just the *bottom row* of the
notation we are using for permutations.

$$\rho_2 = \begin{pmatrix} A & B & C \\ C & A & B \end{pmatrix} = \begin{pmatrix} 1 & 2 & 3 \\ 3 & 1 & 2 \end{pmatrix}$$

```
triangle = SymmetricGroup(3)
rho2 = triangle([3,1,2])
rho2
```

(1,3,2)

```
[rho2(x) for x in triangle.domain()]
```

[3, 1, 2]

The final list comprehension deserves comment. The `.domain()` method gives a lait of the
symbols used for the permutation group `triangle` and then `rho2` is employed with syntax
like it is a function (it *is* a function) to create the images that would occupy the bottom
row.

   With a double list comprehension we can list all six elements of the group in the "bottom
row" format. A good exercise would be to pair up each element with its name as given in
Figure 3.6.

```
[[a(x) for x in triangle.domain()] for a in triangle]
```

[[1, 2, 3], [2, 1, 3], [2, 3, 1], [3, 1, 2], [1, 3, 2], [3, 2, 1]]

Different books, different authors, different software all have different ideas about the order
in which to write multiplication of functions. This textbook builds on the idea of compo-
sition of functions, so that fg is the composition $(fg)(x) = f(g(x))$ and it is natural to
apply g first. Sage takes the opposite view and since we write fg, Sage will understand that
we want to do f first. Neither approach is wrong, and neither is necessarily superior, they
are just different and there are good arguments for either one. When you consult other
books that work with permutation groups, you want to first determine which approach it
takes. (Be aware that this discussion of Sage function composition is limited to permuta-
tions only—"regular" functions in Sage compose in the order you might be familiar with
from a calculus course.)

   The translation here between the text and Sage will be worthwhile practice. Here we
will reprise the discussion at the end of Section 3.1, but reverse the order on each product
to compute Sage-style and exactly mirror what the text does.

```
mu1 = triangle([1,3,2])
mu2 = triangle([3,2,1])
mu3 = triangle([2,1,3])
rho1 = triangle([2,3,1])
product = rho1*mu1
product == mu2
```

True

```
[product(x) for x in triangle.domain()]
```

[3, 2, 1]

```
rho1*mu1 == mu1*rho1
```

False

```
mu1*rho1 == mu3
```

True

Now that we understand that Sage does multiplication in reverse, we can compute the Cayley table for this group. Default behavior is to just name elements of a group as letters, a, b, c, … in the same order that the .list() command would produce the elements of the group. But you can also print the elements in the table as themselves (that uses cycle notation here), or you can give the elements names. We will use u as shorthand for μ and r as shorthand for ρ.

```
triangle.cayley_table()
```

```
*  a b c d e f
 +------------
a| a b c d e f
b| b a f e d c
c| c e d a f b
d| d f a c b e
e| e c b f a d
f| f d e b c a
```

```
triangle.cayley_table(names='elements')
```

```
     *          ()    (1,2) (1,2,3) (1,3,2)    (2,3)    (1,3)
      +-------------------------------------------------------
    ()|          ()    (1,2) (1,2,3) (1,3,2)    (2,3)    (1,3)
  (1,2)|     (1,2)        ()    (1,3)    (2,3) (1,3,2) (1,2,3)
(1,2,3)|  (1,2,3)    (2,3) (1,3,2)        ()    (1,3)    (1,2)
(1,3,2)|  (1,3,2)    (1,3)       () (1,2,3)    (1,2)    (2,3)
  (2,3)|     (2,3) (1,2,3)    (1,2)    (1,3)        () (1,3,2)
  (1,3)|     (1,3) (1,3,2)    (2,3)    (1,2) (1,2,3)        ()
```

```
triangle.cayley_table(names=['id','u3','r1','r2','u1','u2'])
```

```
 *   id u3 r1 r2 u1 u2
  +------------------
id|  id u3 r1 r2 u1 u2
u3|  u3 id u2 u1 r2 r1
r1|  r1 u1 r2 id u2 u3
r2|  r2 u2 id r1 u3 u1
u1|  u1 r1 u3 u2 id r2
u2|  u2 r2 u1 u3 r1 id
```

You should verify that the table above is correct, just like Table 3.2 is correct. Remember that the convention is to multiply a row label times a column label, in that order. However, to do a check across the two tables, you will need to recall the difference in ordering between your textbook and Sage.

Quaternions

Sage implements the quaternions, but the elements are not matrices, but rather are permutations. Despite appearances the structure is identical. It should not matter which version

you have in mind (matrices or permutations) if you build the Cayley table and use the default behavior of using letters to name the elements. As permutations, or as letters, can you identify -1, I, J and K?

```
Q = QuaternionGroup()
[[a(x) for x in Q.domain()] for a in Q]
```

```
[[1, 2, 3, 4, 5, 6, 7, 8], [2, 3, 4, 1, 6, 7, 8, 5],
 [5, 8, 7, 6, 3, 2, 1, 4], [3, 4, 1, 2, 7, 8, 5, 6],
 [6, 5, 8, 7, 4, 3, 2, 1], [8, 7, 6, 5, 2, 1, 4, 3],
 [4, 1, 2, 3, 8, 5, 6, 7], [7, 6, 5, 8, 1, 4, 3, 2]]
```

```
Q.cayley_table()
```

```
*  a b c d e f g h
 +----------------
a| a b c d e f g h
b| b d f g c h a e
c| c e d h g b f a
d| d g h a f e b c
e| e h b f d a c g
f| f c g e a d h b
g| g a e b h c d f
h| h f a c b g e d
```

It should be fairly obvious that a is the identity element of the group (1), either from its behavior in the table, or from its "bottom row" representation as the first element of the list above. And if you prefer, you can ask Sage for a list of its outputs when viewed as a function.

```
id = Q.identity()
[id(x) for x in Q.domain()]
```

```
[1, 2, 3, 4, 5, 6, 7, 8]
```

Now -1 should have the property that $-1 \cdot -1 = 1$. We see that the identity element a is on the diagonal of the Cayley table only when we compute d*d. We can verify this easily, borrowing the fourth "bottom row" element from the list above. With this information, once we locate I, we can easily compute $-I$, and so on.

```
minus_one = Q([3, 4, 1, 2, 7, 8, 5, 6])
minus_one*minus_one == Q.identity()
```

```
True
```

See if you can pair up the letters with all eight elements of the quaternions. Be a bit careful with your names, the symbol I is used by Sage for the imaginary number $i = \sqrt{-1}$ (which we will use below), but Sage will silently let you redefine it to be anything you like. Same goes for using lower-case i in Sage. So call your elements of the quaternions something like QI, QJ, QK to avoid confusion.

As we begin to work with groups it is instructive to work with the actual elements. But many properties of groups are totally independent of the order we use for multiplication, or the names or representations we use for the elements. Here are facts about the quaternions we can compute without any knowledge of just how the elements are written or multiplied.

```
Q.is_finite()
```

```
True
```

```
Q.order()
```

8

```
Q.is_abelian()
```

False

Subgroups

The best techniques for creating subgroups will come in future chapters, but we can create some groups that are naturally subgroups of other groups.

Elements of the quaternions were represented by certain permutations of the integers 1 through 8. We can also build the group of *all* permutations of these eight integers. It gets pretty big, so do not list it unless you want a lot of output! (I dare you.)

```
S8 = SymmetricGroup(8)
a = S8.random_element()
[a(x) for x in S8.domain()]      # random
```

[5, 2, 6, 4, 1, 8, 3, 7]

```
S8.order()
```

40320

The quaternions, Q, is a subgroup of the full group of all permutations, the symmetric group S_8 or S8, and Sage regards this as a property of Q.

```
Q.is_subgroup(S8)
```

True

In Sage the complex numbers are known by the name CC. We can create a list of the elements in the subgroup described in Example 3.16. Then we can verify that this set is a subgroup by examining the Cayley table, using multiplication as the operation.

```
H = [CC(1), CC(-1), CC(I), CC(-I)]
CC.multiplication_table(elements=H,
                        names=['1', '-1', 'i', '-i'])
```

```
*    1 -1  i -i
  +------------
 1|   1 -1  i -i
-1|  -1  1 -i  i
 i|   i -i -1  1
-i|  -i  i  1 -1
```

3.7   Sage Exercises

These exercises are about becoming comfortable working with groups in Sage. Sage worksheets have extensive capabilities for making new cells with carefully formatted text, include support for LATEX syntax to express mathematics. So when a question asks for explanation or commentary, make a new cell and communicate clearly with your audience.**1.**   Create the groups CyclicPermutationGroup(8) and DihedralGroup(4) and name these groups C

and D, respectively. We will understand these constructions better shortly, but for now just understand that both objects you create are actually groups.

2.   Check that C and D have the same size by using the `.order()` method. Determine which group is abelian, and which is not, by using the `.is_abelian()` method.

3.   Use the `.cayley_table()` method to create the Cayley table for each group.

4.   Write a nicely formatted discussion identifying differences between the two groups that are discernible in properties of their Cayley tables. In other words, what is *different* about these two groups that you can "see" in the Cayley tables? (In the Sage notebook, a Shift-click on a blue bar will bring up a mini-word-processor, and you can use use dollar signs to embed mathematics formatted using TEX syntax.)

5.   For C locate the one subgroup of order 4. The group D has three subgroups of order 4. Select one of the three subgroups of D that has a different structure than the subgroup you obtained from C.

The `.subgroups()` method will give you a list of all of the subgroups to help you get started. A Cayley table will help you tell the difference between the two subgroups. What properties of these tables did you use to determine the difference in the structure of the subgroups?

6.   The `.subgroup(elt_list)` method of a group will create the smallest subgroup containing the specified elements of the group, when given the elements as a list `elt_list`. Use this command to discover the shortest list of elements necessary to recreate the subgroups you found in the previous exercise. The equality comparison, `==`, can be used to test if two subgroups are equal.

<div align="center">

$\boxed{4}$

Cyclic Groups

</div>

The groups $\mathbb{Z}$ and $\mathbb{Z}_n$, which are among the most familiar and easily understood groups, are both examples of what are called cyclic groups. In this chapter we will study the properties of cyclic groups and cyclic subgroups, which play a fundamental part in the classification of all abelian groups.

4.1 Cyclic Subgroups

Often a subgroup will depend entirely on a single element of the group; that is, knowing that particular element will allow us to compute any other element in the subgroup.

Example 4.1. Suppose that we consider $3 \in \mathbb{Z}$ and look at all multiples (both positive and negative) of 3. As a set, this is

$$3\mathbb{Z} = \{\ldots, -3, 0, 3, 6, \ldots\}.$$

It is easy to see that $3\mathbb{Z}$ is a subgroup of the integers. This subgroup is completely determined by the element 3 since we can obtain all of the other elements of the group by taking multiples of 3. Every element in the subgroup is "generated" by 3.

Example 4.2. If $H = \{2^n : n \in \mathbb{Z}\}$, then H is a subgroup of the multiplicative group of nonzero rational numbers, $\mathbb{Q}^*$. If $a = 2^m$ and $b = 2^n$ are in H, then $ab^{-1} = 2^m 2^{-n} = 2^{m-n}$ is also in H. By Proposition 3.31, H is a subgroup of $\mathbb{Q}^*$ determined by the element 2.

Theorem 4.3. *Let G be a group and a be any element in G. Then the set*

$$\langle a \rangle = \{a^k : k \in \mathbb{Z}\}$$

is a subgroup of G. Furthermore, $\langle a \rangle$ is the smallest subgroup of G that contains a.

PROOF. The identity is in $\langle a \rangle$ since $a^0 = e$. If g and h are any two elements in $\langle a \rangle$, then by the definition of $\langle a \rangle$ we can write $g = a^m$ and $h = a^n$ for some integers m and n. So $gh = a^m a^n = a^{m+n}$ is again in $\langle a \rangle$. Finally, if $g = a^n$ in $\langle a \rangle$, then the inverse $g^{-1} = a^{-n}$ is also in $\langle a \rangle$. Clearly, any subgroup H of G containing a must contain all the powers of a by closure; hence, H contains $\langle a \rangle$. Therefore, $\langle a \rangle$ is the smallest subgroup of G containing a. $\qquad\square$

Remark 4.4. If we are using the "+" notation, as in the case of the integers under addition, we write $\langle a \rangle = \{na : n \in \mathbb{Z}\}$.

For $a \in G$, we call $\langle a \rangle$ the ***cyclic subgroup*** generated by a. If G contains some element a such that $G = \langle a \rangle$, then G is a ***cyclic group***. In this case a is a ***generator*** of G. If a is an element of a group G, we define the ***order*** of a to be the smallest positive integer n such that $a^n = e$, and we write $|a| = n$. If there is no such integer n, we say that the order of a is infinite and write $|a| = \infty$ to denote the order of a.

<div align="center">

59

</div>

Example 4.5. Notice that a cyclic group can have more than a single generator. Both 1 and 5 generate $\mathbb{Z}_6$; hence, $\mathbb{Z}_6$ is a cyclic group. Not every element in a cyclic group is necessarily a generator of the group. The order of $2 \in \mathbb{Z}_6$ is 3. The cyclic subgroup generated by 2 is $\langle 2 \rangle = \{0, 2, 4\}$.

The groups $\mathbb{Z}$ and $\mathbb{Z}_n$ are cyclic groups. The elements 1 and -1 are generators for $\mathbb{Z}$. We can certainly generate $\mathbb{Z}_n$ with 1 although there may be other generators of $\mathbb{Z}_n$, as in the case of $\mathbb{Z}_6$.

Example 4.6. The group of units, $U(9)$, in $\mathbb{Z}_9$ is a cyclic group. As a set, $U(9)$ is $\{1, 2, 4, 5, 7, 8\}$. The element 2 is a generator for $U(9)$ since

$$2^1 = 2 \qquad 2^2 = 4$$
$$2^3 = 8 \qquad 2^4 = 7$$
$$2^5 = 5 \qquad 2^6 = 1.$$

Example 4.7. Not every group is a cyclic group. Consider the symmetry group of an equilateral triangle S_3. The multiplication table for this group is Table 3.7. The subgroups of S_3 are shown in Figure 4.8. Notice that every subgroup is cyclic; however, no single element generates the entire group.
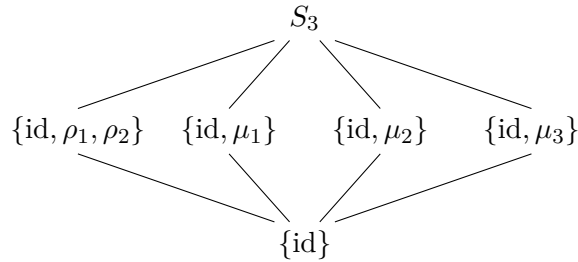
Figure 4.8: Subgroups of S_3

Theorem 4.9. *Every cyclic group is abelian.*

PROOF.   Let G be a cyclic group and $a \in G$ be a generator for G. If g and h are in G, then they can be written as powers of a, say $g = a^r$ and $h = a^s$. Since

$$gh = a^r a^s = a^{r+s} = a^{s+r} = a^s a^r = hg,$$

G is abelian.   □

Subgroups of Cyclic Groups

We can ask some interesting questions about cyclic subgroups of a group and subgroups of a cyclic group. If G is a group, which subgroups of G are cyclic? If G is a cyclic group, what type of subgroups does G possess?

Theorem 4.10. *Every subgroup of a cyclic group is cyclic.*

PROOF.   The main tools used in this proof are the division algorithm and the Principle of Well-Ordering. Let G be a cyclic group generated by a and suppose that H is a subgroup of G. If $H = \{e\}$, then trivially H is cyclic. Suppose that H contains some other element g distinct from the identity. Then g can be written as a^n for some integer n. Since H is a

subgroup, $g^{-1} = a^{-n}$ must also be in H. Since either n or $-n$ is positive, we can assume that H contains positive powers of a and $n > 0$. Let m be the smallest natural number such that $a^m \in H$. Such an m exists by the Principle of Well-Ordering.

We claim that $h = a^m$ is a generator for H. We must show that every $h' \in H$ can be written as a power of h. Since $h' \in H$ and H is a subgroup of G, $h' = a^k$ for some integer k. Using the division algorithm, we can find numbers q and r such that $k = mq + r$ where $0 \le r < m$; hence,

$$a^k = a^{mq+r} = (a^m)^q a^r = h^q a^r.$$

So $a^r = a^k h^{-q}$. Since a^k and h^{-q} are in H, a^r must also be in H. However, m was the smallest positive number such that a^m was in H; consequently, $r = 0$ and so $k = mq$. Therefore,

$$h' = a^k = a^{mq} = h^q$$

and H is generated by h. $\square$

Corollary 4.11. *The subgroups of $\mathbb{Z}$ are exactly $n\mathbb{Z}$ for $n = 0, 1, 2, \ldots$.*

Proposition 4.12. *Let G be a cyclic group of order n and suppose that a is a generator for G. Then $a^k = e$ if and only if n divides k.*

PROOF. First suppose that $a^k = e$. By the division algorithm, $k = nq + r$ where $0 \le r < n$; hence,

$$e = a^k = a^{nq+r} = a^{nq} a^r = ea^r = a^r.$$

Since the smallest positive integer m such that $a^m = e$ is n, $r = 0$.

Conversely, if n divides k, then $k = ns$ for some integer s. Consequently,

$$a^k = a^{ns} = (a^n)^s = e^s = e. \qquad \square$$

Theorem 4.13. *Let G be a cyclic group of order n and suppose that $a \in G$ is a generator of the group. If $b = a^k$, then the order of b is n/d, where $d = \gcd(k, n)$.*

PROOF. We wish to find the smallest integer m such that $e = b^m = a^{km}$. By Proposition 4.12, this is the smallest integer m such that n divides km or, equivalently, n/d divides $m(k/d)$. Since d is the greatest common divisor of n and k, n/d and k/d are relatively prime. Hence, for n/d to divide $m(k/d)$ it must divide m. The smallest such m is n/d. $\square$

Corollary 4.14. *The generators of $\mathbb{Z}_n$ are the integers r such that $1 \le r < n$ and $\gcd(r, n) = 1$.*

Example 4.15. Let us examine the group $\mathbb{Z}_{16}$. The numbers 1, 3, 5, 7, 9, 11, 13, and 15 are the elements of $\mathbb{Z}_{16}$ that are relatively prime to 16. Each of these elements generates $\mathbb{Z}_{16}$. For example,

$$
\begin{array}{lll}
1 \cdot 9 = 9 & 2 \cdot 9 = 2 & 3 \cdot 9 = 11 \\
4 \cdot 9 = 4 & 5 \cdot 9 = 13 & 6 \cdot 9 = 6 \\
7 \cdot 9 = 15 & 8 \cdot 9 = 8 & 9 \cdot 9 = 1 \\
10 \cdot 9 = 10 & 11 \cdot 9 = 3 & 12 \cdot 9 = 12 \\
13 \cdot 9 = 5 & 14 \cdot 9 = 14 & 15 \cdot 9 = 7.
\end{array}
$$

4.2   Multiplicative Group of Complex Numbers

The **complex numbers** are defined as

$$\mathbb{C} = \{a + bi : a, b \in \mathbb{R}\},$$

where $i^2 = -1$. If $z = a + bi$, then a is the **real part** of z and b is the **imaginary part** of z.

To add two complex numbers $z = a + bi$ and $w = c + di$, we just add the corresponding real and imaginary parts:

$$z + w = (a + bi) + (c + di) = (a + c) + (b + d)i.$$

Remembering that $i^2 = -1$, we multiply complex numbers just like polynomials. The product of z and w is

$$(a + bi)(c + di) = ac + bdi^2 + adi + bci = (ac - bd) + (ad + bc)i.$$

Every nonzero complex number $z = a + bi$ has a multiplicative inverse; that is, there exists a $z^{-1} \in \mathbb{C}^*$ such that $zz^{-1} = z^{-1}z = 1$. If $z = a + bi$, then

$$z^{-1} = \frac{a - bi}{a^2 + b^2}.$$

The **complex conjugate** of a complex number $z = a + bi$ is defined to be $\overline{z} = a - bi$. The **absolute value** or **modulus** of $z = a + bi$ is $|z| = \sqrt{a^2 + b^2}$.

Example 4.16. Let $z = 2 + 3i$ and $w = 1 - 2i$. Then

$$z + w = (2 + 3i) + (1 - 2i) = 3 + i$$

and

$$zw = (2 + 3i)(1 - 2i) = 8 - i.$$

Also,

$$z^{-1} = \frac{2}{13} - \frac{3}{13}i$$
$$|z| = \sqrt{13}$$
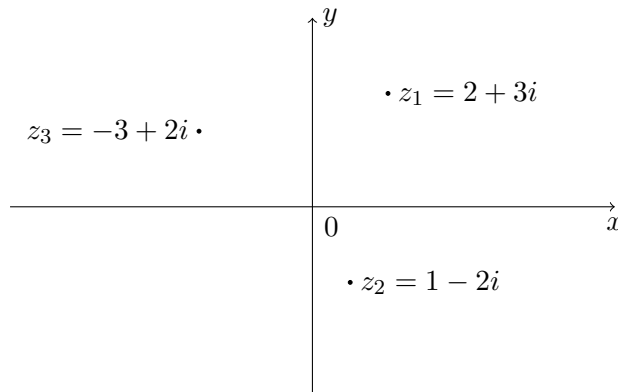$$\overline{z} = 2 - 3i.$$

Figure 4.17: Rectangular coordinates of a complex number

There are several ways of graphically representing complex numbers. We can represent a complex number $z = a + bi$ as an ordered pair on the xy plane where a is the x (or real) coordinate and b is the y (or imaginary) coordinate. This is called the **rectangular** or **Cartesian** representation. The rectangular representations of $z_1 = 2 + 3i$, $z_2 = 1 - 2i$, and $z_3 = -3 + 2i$ are depicted in Figure 4.17.
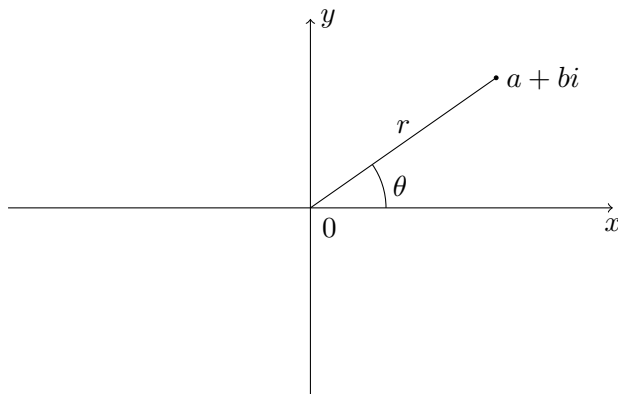
Figure 4.18: Polar coordinates of a complex number

Nonzero complex numbers can also be represented using **polar coordinates**. To specify any nonzero point on the plane, it suffices to give an angle θ from the positive x axis in the counterclockwise direction and a distance r from the origin, as in Figure 4.18. We can see that

$$z = a + bi = r(\cos\theta + i\sin\theta).$$

Hence,

$$r = |z| = \sqrt{a^2 + b^2}$$

and

$$a = r\cos\theta$$
$$b = r\sin\theta.$$

We sometimes abbreviate $r(\cos\theta + i\sin\theta)$ as $r\operatorname{cis}\theta$. To assure that the representation of z is well-defined, we also require that $0° \le \theta < 360°$. If the measurement is in radians, then $0 \le \theta < 2\pi$.

Example 4.19. Suppose that $z = 2\operatorname{cis} 60°$. Then

$$a = 2\cos 60° = 1$$

and

$$b = 2\sin 60° = \sqrt{3}.$$

Hence, the rectangular representation is $z = 1 + \sqrt{3}\, i$.

Conversely, if we are given a rectangular representation of a complex number, it is often useful to know the number's polar representation. If $z = 3\sqrt{2} - 3\sqrt{2}\, i$, then

$$r = \sqrt{a^2 + b^2} = \sqrt{36} = 6$$

and

$$\theta = \arctan\left(\frac{b}{a}\right) = \arctan(-1) = 315°,$$

so $3\sqrt{2} - 3\sqrt{2}\, i = 6\operatorname{cis} 315°$.

The polar representation of a complex number makes it easy to find products and powers of complex numbers. The proof of the following proposition is straightforward and is left as an exercise.

Proposition 4.20. *Let $z = r \operatorname{cis} \theta$ and $w = s \operatorname{cis} \phi$ be two nonzero complex numbers. Then*

$$zw = rs \operatorname{cis}(\theta + \phi).$$

Example 4.21. If $z = 3 \operatorname{cis}(\pi/3)$ and $w = 2 \operatorname{cis}(\pi/6)$, then $zw = 6 \operatorname{cis}(\pi/2) = 6i$.

Theorem 4.22 DeMoivre. *Let $z = r \operatorname{cis} \theta$ be a nonzero complex number. Then*

$$[r \operatorname{cis} \theta]^n = r^n \operatorname{cis}(n\theta)$$

for $n = 1, 2, \ldots$.

PROOF.   We will use induction on n. For $n = 1$ the theorem is trivial. Assume that the theorem is true for all k such that $1 \leq k \leq n$. Then

$$
\begin{aligned}
z^{n+1} &= z^n z \\
&= r^n(\cos n\theta + i \sin n\theta) r(\cos \theta + i \sin \theta) \\
&= r^{n+1}[(\cos n\theta \cos \theta - \sin n\theta \sin \theta) + i(\sin n\theta \cos \theta + \cos n\theta \sin \theta)] \\
&= r^{n+1}[\cos(n\theta + \theta) + i \sin(n\theta + \theta)] \\
&= r^{n+1}[\cos(n + 1)\theta + i \sin(n + 1)\theta]. \qquad \square
\end{aligned}
$$

Example 4.23. Suppose that $z = 1 + i$ and we wish to compute z^{10}. Rather than computing $(1 + i)^{10}$ directly, it is much easier to switch to polar coordinates and calculate z^{10} using DeMoivre's Theorem:

$$
\begin{aligned}
z^{10} &= (1 + i)^{10} \\
&= \left(\sqrt{2} \operatorname{cis}\left(\frac{\pi}{4}\right)\right)^{10} \\
&= (\sqrt{2})^{10} \operatorname{cis}\left(\frac{5\pi}{2}\right) \\
&= 32 \operatorname{cis}\left(\frac{\pi}{2}\right) \\
&= 32i.
\end{aligned}
$$

The Circle Group and the Roots of Unity

The multiplicative group of the complex numbers, $\mathbb{C}^*$, possesses some interesting subgroups. Whereas $\mathbb{Q}^*$ and $\mathbb{R}^*$ have no interesting subgroups of finite order, $\mathbb{C}^*$ has many. We first consider the ***circle group***,

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

The following proposition is a direct result of Proposition 4.20.

Proposition 4.24. *The circle group is a subgroup of $\mathbb{C}^*$.*

Although the circle group has infinite order, it has many interesting finite subgroups. Suppose that $H = \{1, -1, i, -i\}$. Then H is a subgroup of the circle group. Also, $1, -1, i$, and $-i$ are exactly those complex numbers that satisfy the equation $z^4 = 1$. The complex numbers satisfying the equation $z^n = 1$ are called the ***nth roots of unity***.

Theorem 4.25. *If $z^n = 1$, then the nth roots of unity are*

$$z = \operatorname{cis}\left(\frac{2k\pi}{n}\right),$$

where $k = 0, 1, \ldots, n-1$. Furthermore, the nth roots of unity form a cyclic subgroup of $\mathbb{T}$ of order n

PROOF. By DeMoivre's Theorem,

$$z^n = \operatorname{cis}\left(n\frac{2k\pi}{n}\right) = \operatorname{cis}(2k\pi) = 1.$$

The z's are distinct since the numbers $2k\pi/n$ are all distinct and are greater than or equal to 0 but less than 2π. The fact that these are all of the roots of the equation $z^n = 1$ follows from from Corollary 17.9, which states that a polynomial of degree n can have at most n roots. We will leave the proof that the nth roots of unity form a cyclic subgroup of $\mathbb{T}$ as an exercise. $\qquad\square$

A generator for the group of the nth roots of unity is called a **primitive nth root of unity**.

Example 4.26. The 8th roots of unity can be represented as eight equally spaced points on the unit circle (Figure 4.27). The primitive 8th roots of unity are

$$\omega = \frac{\sqrt{2}}{2} + \frac{\sqrt{2}}{2}i$$

$$\omega^3 = -\frac{\sqrt{2}}{2} + \frac{\sqrt{2}}{2}i$$

$$\omega^5 = -\frac{\sqrt{2}}{2} - \frac{\sqrt{2}}{2}i$$

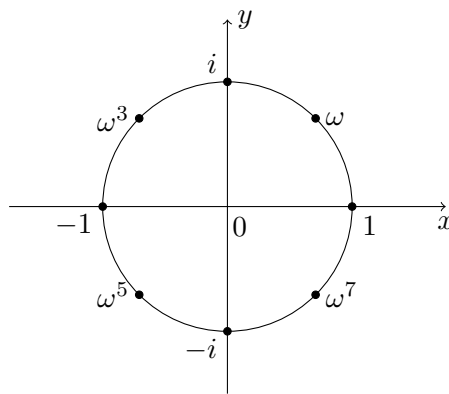$$\omega^7 = \frac{\sqrt{2}}{2} - \frac{\sqrt{2}}{2}i.$$

Figure 4.27: 8th roots of unity

4.3 The Method of Repeated Squares

Computing large powers can be very time-consuming. Just as anyone can compute 2^2 or 2^8, everyone knows how to compute

$$2^{2^{1,000,000}}.$$

However, such numbers are so large that we do not want to attempt the calculations; moreover, past a certain point the computations would not be feasible even if we had every computer in the world at our disposal. Even writing down the decimal representation of a very large number may not be reasonable. It could be thousands or even millions of digits long. However, if we could compute something like

$$2^{37,398,332} \pmod{46,389},$$

we could very easily write the result down since it would be a number between 0 and 46,388. If we want to compute powers modulo n quickly and efficiently, we will have to be clever.[1]

The first thing to notice is that any number a can be written as the sum of distinct powers of 2; that is, we can write

$$a = 2^{k_1} + 2^{k_2} + \cdots + 2^{k_n},$$

where $k_1 < k_2 < \cdots < k_n$. This is just the binary representation of a. For example, the binary representation of 57 is 111001, since we can write $57 = 2^0 + 2^3 + 2^4 + 2^5$.

The laws of exponents still work in $\mathbb{Z}_n$; that is, if $b \equiv a^x \pmod{n}$ and $c \equiv a^y \pmod{n}$, then $bc \equiv a^{x+y} \pmod{n}$. We can compute $a^{2^k} \pmod{n}$ in k multiplications by computing

$$a^{2^0} \pmod{n}$$
$$a^{2^1} \pmod{n}$$
$$\vdots$$
$$a^{2^k} \pmod{n}.$$

Each step involves squaring the answer obtained in the previous step, dividing by n, and taking the remainder.

Example 4.28. We will compute $271^{321} \pmod{481}$. Notice that

$$321 = 2^0 + 2^6 + 2^8;$$

hence, computing $271^{321} \pmod{481}$ is the same as computing

$$271^{2^0 + 2^6 + 2^8} \equiv 271^{2^0} \cdot 271^{2^6} \cdot 271^{2^8} \pmod{481}.$$

So it will suffice to compute $271^{2^i} \pmod{481}$ where $i = 0, 6, 8$. It is very easy to see that

$$271^{2^1} = 73,441 \equiv 329 \pmod{481}.$$

We can square this result to obtain a value for $271^{2^2} \pmod{481}$:

$$271^{2^2} \equiv (271^{2^1})^2 \pmod{481}$$
$$\equiv (329)^2 \pmod{481}$$
$$\equiv 108,241 \pmod{481}$$
$$\equiv 16 \pmod{481}.$$

We are using the fact that $(a^{2^n})^2 \equiv a^{2 \cdot 2^n} \equiv a^{2^{n+1}} \pmod{n}$. Continuing, we can calculate

$$271^{2^6} \equiv 419 \pmod{481}$$

---

[1]The results in this section are needed only in Chapter 7

and
$$271^{2^8} \equiv 16 \pmod{481}.$$

Therefore,

$$
\begin{aligned}
271^{321} &\equiv 271^{2^0 + 2^6 + 2^8} \pmod{481} \\
&\equiv 271^{2^0} \cdot 271^{2^6} \cdot 271^{2^8} \pmod{481} \\
&\equiv 271 \cdot 419 \cdot 16 \pmod{481} \\
&\equiv 1{,}816{,}784 \pmod{481} \\
&\equiv 47 \pmod{481}.
\end{aligned}
$$

The method of repeated squares will prove to be a very useful tool when we explore RSA cryptography in Chapter 7. To encode and decode messages in a reasonable manner under this scheme, it is necessary to be able to quickly compute large powers of integers mod n.

4.4 Exercises

1. Prove or disprove each of the following statements.

(a) All of the generators of $\mathbb{Z}_{60}$ are prime.

(b) $U(8)$ is cyclic.

(c) $\mathbb{Q}$ is cyclic.

(d) If every proper subgroup of a group G is cyclic, then G is a cyclic group.

(e) A group with a finite number of subgroups is finite.

2. Find the order of each of the following elements.

(a) $5 \in \mathbb{Z}_{12}$

(b) $\sqrt{3} \in \mathbb{R}$

(c) $\sqrt{3} \in \mathbb{R}^*$

(d) $-i \in \mathbb{C}^*$

(e) $72 \in \mathbb{Z}_{240}$

(f) $312 \in \mathbb{Z}_{471}$

3. List all of the elements in each of the following subgroups.

(a) The subgroup of $\mathbb{Z}$ generated by 7

(b) The subgroup of $\mathbb{Z}_{24}$ generated by 15

(c) All subgroups of $\mathbb{Z}_{12}$

(d) All subgroups of $\mathbb{Z}_{60}$

(e) All subgroups of $\mathbb{Z}_{13}$

(f) All subgroups of $\mathbb{Z}_{48}$

(g) The subgroup generated by 3 in $U(20)$

(h) The subgroup generated by 5 in $U(18)$

(i) The subgroup of $\mathbb{R}^*$ generated by 7

(j) The subgroup of $\mathbb{C}^*$ generated by i where $i^2 = -1$

(k) The subgroup of $\mathbb{C}^*$ generated by $2i$

(l) The subgroup of $\mathbb{C}^*$ generated by $(1 + i)/\sqrt{2}$

(m) The subgroup of $\mathbb{C}^*$ generated by $(1 + \sqrt{3}\, i)/2$

4.   Find the subgroups of $GL_2(\mathbb{R})$ generated by each of the following matrices.

(a) $\begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}$
  (c) $\begin{pmatrix} 1 & -1 \\ 1 & 0 \end{pmatrix}$
  (e) $\begin{pmatrix} 1 & -1 \\ -1 & 0 \end{pmatrix}$

(b) $\begin{pmatrix} 0 & 1/3 \\ 3 & 0 \end{pmatrix}$
  (d) $\begin{pmatrix} 1 & -1 \\ 0 & 1 \end{pmatrix}$
  (f) $\begin{pmatrix} \sqrt{3}/2 & 1/2 \\ -1/2 & \sqrt{3}/2 \end{pmatrix}$

5.   Find the order of every element in $\mathbb{Z}_{18}$.

6.   Find the order of every element in the symmetry group of the square, D_4.

7.   What are all of the cyclic subgroups of the quaternion group, Q_8?

8.   List all of the cyclic subgroups of $U(30)$.

9.   List every generator of each subgroup of order 8 in $\mathbb{Z}_{32}$.

10.   Find all elements of finite order in each of the following groups. Here the "$*$" indicates the set with zero removed.

(a) $\mathbb{Z}$
  (b) $\mathbb{Q}^*$
  (c) $\mathbb{R}^*$

11.   If $a^{24} = e$ in a group G, what are the possible orders of a?

12.   Find a cyclic group with exactly one generator. Can you find cyclic groups with exactly two generators? Four generators? How about n generators?

13.   For $n \le 20$, which groups $U(n)$ are cyclic? Make a conjecture as to what is true in general. Can you prove your conjecture?

14.   Let

$$A = \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix} \qquad \text{and} \qquad B = \begin{pmatrix} 0 & -1 \\ 1 & -1 \end{pmatrix}$$

be elements in $GL_2(\mathbb{R})$. Show that A and B have finite orders but AB does not.

15.   Evaluate each of the following.

(a) $(3 - 2i) + (5i - 6)$
  (d) $(9 - i)\overline{(9 - i)}$

(b) $(4 - 5i) - \overline{(4i - 4)}$
  (e) i^{45}

(c) $(5 - 4i)(7 + 2i)$
  (f) $(1 + i) + \overline{(1 + i)}$

16.   Convert the following complex numbers to the form $a + bi$.

(a) $2\operatorname{cis}(\pi/6)$
  (c) $3\operatorname{cis}(\pi)$

(b) $5\operatorname{cis}(9\pi/4)$
  (d) $\operatorname{cis}(7\pi/4)/2$

17.   Change the following complex numbers to polar representation.

(a) $1 - i$ (c) $2 + 2i$ (e) $-3i$

(b) -5 (d) $\sqrt{3} + i$ (f) $2i + 2\sqrt{3}$

18. Calculate each of the following expressions.

(a) $(1 + i)^{-1}$

(b) $(1 - i)^6$

(c) $(\sqrt{3} + i)^5$

(d) $(-i)^{10}$

(e) $((1 - i)/2)^4$

(f) $(-\sqrt{2} - \sqrt{2}\,i)^{12}$

(g) $(-2 + 2i)^{-5}$

19. Prove each of the following statements.

(a) $|z| = |\bar{z}|$

(b) $z\bar{z} = |z|^2$

(c) $z^{-1} = \bar{z}/|z|^2$

(d) $|z + w| \leq |z| + |w|$

(e) $|z - w| \geq ||z| - |w||$

(f) $|zw| = |z||w|$

20. List and graph the 6th roots of unity. What are the generators of this group? What are the primitive 6th roots of unity?

21. List and graph the 5th roots of unity. What are the generators of this group? What are the primitive 5th roots of unity?

22. Calculate each of the following.

(a) $292^{3171} \pmod{582}$

(b) $2557^{341} \pmod{5681}$

(c) $2071^{9521} \pmod{4724}$

(d) $971^{321} \pmod{765}$

23. Let $a, b \in G$. Prove the following statements.

(a) The order of a is the same as the order of a^{-1}.

(b) For all $g \in G$, $|a| = |g^{-1}ag|$.

(c) The order of ab is the same as the order of ba.

24. Let p and q be distinct primes. How many generators does $\mathbb{Z}_{pq}$ have?

25. Let p be prime and r be a positive integer. How many generators does $\mathbb{Z}_{p^r}$ have?

26. Prove that $\mathbb{Z}_p$ has no nontrivial subgroups if p is prime.

27. If g and h have orders 15 and 16 respectively in a group G, what is the order of $\langle g \rangle \cap \langle h \rangle$?

28. Let a be an element in a group G. What is a generator for the subgroup $\langle a^m \rangle \cap \langle a^n \rangle$?

29. Prove that $\mathbb{Z}_n$ has an even number of generators for $n > 2$.

30. Suppose that G is a group and let a, $b \in G$. Prove that if $|a| = m$ and $|b| = n$ with $\gcd(m, n) = 1$, then $\langle a \rangle \cap \langle b \rangle = \{e\}$.

31. Let G be an abelian group. Show that the elements of finite order in G form a subgroup. This subgroup is called the ***torsion subgroup*** of G.

32. Let G be a finite cyclic group of order n generated by x. Show that if $y = x^k$ where $\gcd(k, n) = 1$, then y must be a generator of G.

33.   If G is an abelian group that contains a pair of cyclic subgroups of order 2, show that G must contain a subgroup of order 4. Does this subgroup have to be cyclic?

34.   Let G be an abelian group of order pq where $\gcd(p,q) = 1$. If G contains elements a and b of order p and q respectively, then show that G is cyclic.

35.   Prove that the subgroups of $\mathbb{Z}$ are exactly $n\mathbb{Z}$ for $n = 0, 1, 2, \ldots$.

36.   Prove that the generators of $\mathbb{Z}_n$ are the integers r such that $1 \leq r < n$ and $\gcd(r, n) = 1$.

37.   Prove that if G has no proper nontrivial subgroups, then G is a cyclic group.

38.   Prove that the order of an element in a cyclic group G must divide the order of the group.

39.   Prove that if G is a cyclic group of order m and $d \mid m$, then G must have a subgroup of order d.

40.   For what integers n is -1 an nth root of unity?

41.   If $z = r(\cos\theta + i\sin\theta)$ and $w = s(\cos\phi + i\sin\phi)$ are two nonzero complex numbers, show that

$$zw = rs[\cos(\theta + \phi) + i\sin(\theta + \phi)].$$

42.   Prove that the circle group is a subgroup of $\mathbb{C}^*$.

43.   Prove that the nth roots of unity form a cyclic subgroup of $\mathbb{T}$ of order n.

44.   Let $\alpha \in \mathbb{T}$. Prove that $\alpha^m = 1$ and $\alpha^n = 1$ if and only if $\alpha^d = 1$ for $d = \gcd(m, n)$.

45.   Let $z \in \mathbb{C}^*$. If $|z| \neq 1$, prove that the order of z is infinite.

46.   Let $z = \cos\theta + i\sin\theta$ be in $\mathbb{T}$ where $\theta \in \mathbb{Q}$. Prove that the order of z is infinite.

4.5   Programming Exercises

1.   Write a computer program that will write any decimal number as the sum of distinct powers of 2. What is the largest integer that your program will handle?

2.   Write a computer program to calculate $a^x \pmod{n}$ by the method of repeated squares. What are the largest values of n and x that your program will accept?

References and Suggested Readings

[1]   Koblitz, N. *A Course in Number Theory and Cryptography.* 2nd ed. Springer, New York, 1994.

[2]   Pomerance, C. "Cryptology and Computational Number Theory—An Introduction," in *Cryptology and Computational Number Theory*, Pomerance, C., ed. Proceedings of Symposia in Applied Mathematics, vol. 42, American Mathematical Society, Providence, RI, 1990. Thisbook gives an excellent account of how the method of repeated squares is used in cryptography.

4.6 Sage

Cyclic groups are very important, so it is no surprise that they appear in many different forms in Sage. Each is slightly different, and no one implementation is ideal for an introduction, but together they can illustrate most of the important ideas. Here is a guide to the various ways to construct, and study, a cyclic group in Sage.

Infinite Cyclic Groups

In Sage, the integers $\mathbb{Z}$ are constructed with ZZ. To build the infinite cyclic group such as $3\mathbb{Z}$ from Example 4.1, simply use 3*ZZ. As an infinite set, there is not a whole lot you can do with this. You can test if integers are in this set, or not. You can also recall the generator with the .gen() command.

```
G = 3*ZZ
-12 in G
```

True

```
37 in G
```

False

```
G.gen()
```

3

Additive Cyclic Groups

The additive cyclic group $\mathbb{Z}_n$ can be built as a special case of a more general Sage construction. First we build $\mathbb{Z}_{14}$ and capture its generator. Throughout, pay close attention to the use of parentheses and square brackets for when you experiment on your own.

```
G = AdditiveAbelianGroup([14])
G.order()
```

14

```
G.list()
```

```
[(0), (1), (2), (3), (4), (5), (6), (7),
 (8), (9), (10), (11), (12), (13)]
```

```
a = G.gen(0)
a
```

(1)

You can compute in this group, by using the generator, or by using new elements formed by coercing integers into the group, or by taking the result of operations on other elements. And we can compute the order of elements in this group. Notice that we can perform repeated additions with the shortcut of taking integer multiples of an element.

```
a + a
```

(2)

```
a + a + a + a
```

(4)

```
4*a
```

(4)

```
37*a
```

(9)

We can create, and then compute with, new elements of the group by coercing an integer (in a list of length 1) into the group. You may get a `DeprecationWarning` the first time you use this syntax to create a new element. The mysterious warning can be safely ignored.

```
G([2])
```

```
doctest:...: DeprecationWarning: The default behaviour changed! If you
*really* want a linear combination of smith generators, use
.linear_combination_of_smith_form_gens.
See http://trac.sagemath.org/16261 for details.
(2)
```

```
b = G([2]); b
```

(2)

```
b + b
```

(4)

```
2*b == 4*a
```

True

```
7*b
```

(0)

```
b.order()
```

7

```
c = a - 6*b; c
```

(3)

```
c + c + c + c
```

(12)

```
c.order()
```

14

It is possible to create cyclic subgroups, from an element designated to be the new generator. Unfortunately, to do this requires the `.submodule()` method (which should be renamed in Sage).

```
H = G.submodule([b]); H
```

Additive abelian group isomorphic to Z/7

```
H.list()
```

[(0), (2), (4), (6), (8), (10), (12)]

```
H.order()
```

7

```
e = H.gen(0); e
```

(2)

```
3*e
```

(6)

```
e.order()
```

7

The cyclic subgroup `H` just created has more than one generator. We can test this by building a new subgroup and comparing the two subgroups.

```
f = 12*a; f
```

(12)

```
f.order()
```

7

```
K = G.submodule([f]); K
```

Additive abelian group isomorphic to Z/7

```
K.order()
```

7

```
K.list()
```

[(0), (2), (4), (6), (8), (10), (12)]

```
K.gen(0)
```

(2)

```
H == K
```

```
True
```

Certainly the list of elements, and the common generator of (2) lead us to belive that H
and K are the same, but the comparison in the last line leaves no doubt.

Results in this section, especially Theorem 4.13 and Corollary 4.14, can be investigated
by creating generators of subgroups from a generator of one additive cyclic group, creating
the subgroups, and computing the orders of both elements and orders of groups.

Abstract Multiplicative Cyclic Groups

We can create an abstract cyclic group in the style of Theorems 4.3, 4.9, 4.10. In the
syntax below a is a name for the generator, and 14 is the order of the element. Notice that
the notation is now multiplicative, so we multiply elements, and repeated products can be
written as powers.

```
G.<a> = AbelianGroup([14])
G.order()
```

```
14
```

```
G.list()
```

```
(1, a, a^2, a^3, a^4, a^5, a^6, a^7, a^8, a^9, a^10, a^11, a^12, a^13)
```

```
a.order()
```

```
14
```

Computations in the group are similar to before, only with different notation. Now products,
with repeated products written as exponentiation.

```
b = a^2
b.order()
```

```
7
```

```
b*b*b
```

```
a^6
```

```
c = a^7
c.order()
```

```
2
```

```
c^2
```

```
1
```

```
b*c
```

```
a^9
```

```
b^37*c^42
```

```
a^4
```

Subgroups can be formed with a `.subgroup()` command. But do not try to list the contents of a subgroup, it'll look strangely unfamiliar. Also, comparison of subgroups is not implemented.

```
H = G.subgroup([a^2])
H.order()
```

7

```
K = G.subgroup([a^12])
K.order()
```

7

```
L = G.subgroup([a^4])
H == L
```

False

One advantage of this implementation is the possibility to create all possible subgroups. Here we create the list of subgroups, extract one in particular (the third), and check its order.

```
allsg = G.subgroups(); allsg
```

```
[Multiplicative Abelian subgroup isomorphic to C2 x C7 generated by {a},
 Multiplicative Abelian subgroup isomorphic to C7 generated by {a^2},
 Multiplicative Abelian subgroup isomorphic to C2 generated by {a^7},
 Trivial Abelian subgroup]
```

```
sub = allsg[2]
sub.order()
```

2

Cyclic Permutation Groups

We will learn more about permutation groups in the next chapter. But we will mention here that it is easy to create cyclic groups as permutation groups, and a variety of methods are available for working with them, even if the actual elements get a bit cumbersome to work with. As before, notice that the notation and syntax is multiplicative.

```
G=CyclicPermutationGroup(14)
a = G.gen(0); a
```

(1,2,3,4,5,6,7,8,9,10,11,12,13,14)

```
b = a^2
b = a^2; b
```

(1,3,5,7,9,11,13)(2,4,6,8,10,12,14)

```
b.order()
```

7

```
a*a*b*b*b
```

(1,9,3,11,5,13,7)(2,10,4,12,6,14,8)

```
c = a^37*b^26; c
```

(1,6,11,2,7,12,3,8,13,4,9,14,5,10)

```
c.order()
```

14

We can create subgroups, check their orders, and list their elements.

```
H = G.subgroup([a^2])
H.order()
```

7

```
H.gen(0)
```

(1,3,5,7,9,11,13)(2,4,6,8,10,12,14)

```
H.list()
```

```
[(),
 (1,3,5,7,9,11,13)(2,4,6,8,10,12,14),
 (1,5,9,13,3,7,11)(2,6,10,14,4,8,12),
 (1,7,13,5,11,3,9)(2,8,14,6,12,4,10),
 (1,9,3,11,5,13,7)(2,10,4,12,6,14,8),
 (1,11,7,3,13,9,5)(2,12,8,4,14,10,6),
 (1,13,11,9,7,5,3)(2,14,12,10,8,6,4)]
```

It could help to visualize this group, and the subgroup, as rotations of a regular 12-gon with the vertices labeled with the integers 1 through 12. This is not the full group of symmetries, since it does not include reflections, just the 12 rotations.

Cayley Tables

As groups, each of the examples above (groups and subgroups) have Cayley tables implemented. Since the groups are cyclic, and their subgroups are therefore cyclic, the Cayley tables should have a similar "cyclic" pattern. Note that the letters used in the default table are generic, and are not related to the letters used above for specific elements — they just match up with the group elements in the order given by `.list()`.

```
G.<a> = AbelianGroup([14])
G.cayley_table()
```

```
*   a b c d e f g h i j k l m n
 +---------------------------
a|  a b c d e f g h i j k l m n
b|  b c d e f g h i j k l m n a
c|  c d e f g h i j k l m n a b
d|  d e f g h i j k l m n a b c
e|  e f g h i j k l m n a b c d
f|  f g h i j k l m n a b c d e
g|  g h i j k l m n a b c d e f
h|  h i j k l m n a b c d e f g
i|  i j k l m n a b c d e f g h
```

```
j|  j k l m n a b c d e f g h i
k|  k l m n a b c d e f g h i j
l|  l m n a b c d e f g h i j k
m|  m n a b c d e f g h i j k l
n|  n a b c d e f g h i j k l m
```

If the real names of the elements are not too complicated, the table could be more informative using these names.

```
K.<b> = AbelianGroup([10])
K.cayley_table(names='elements')
```

```
  *     1    b b^2 b^3 b^4 b^5 b^6 b^7 b^8 b^9
   +---------------------------------------
  1|    1    b b^2 b^3 b^4 b^5 b^6 b^7 b^8 b^9
  b|    b  b^2 b^3 b^4 b^5 b^6 b^7 b^8 b^9    1
b^2|  b^2 b^3 b^4 b^5 b^6 b^7 b^8 b^9    1    b
b^3|  b^3 b^4 b^5 b^6 b^7 b^8 b^9    1    b b^2
b^4|  b^4 b^5 b^6 b^7 b^8 b^9    1    b b^2 b^3
b^5|  b^5 b^6 b^7 b^8 b^9    1    b b^2 b^3 b^4
b^6|  b^6 b^7 b^8 b^9    1    b b^2 b^3 b^4 b^5
b^7|  b^7 b^8 b^9    1    b b^2 b^3 b^4 b^5 b^6
b^8|  b^8 b^9    1    b b^2 b^3 b^4 b^5 b^6 b^7
b^9|  b^9    1    b b^2 b^3 b^4 b^5 b^6 b^7 b^8
```

Complex Roots of Unity

The finite cyclic subgroups of $\mathbb{T}$, generated by a primitive nth root of unity are implemented as a more general construction in Sage, known as a cyclotomic field. If you concentrate on just the multiplication of powers of a generator (and ignore the infinitely many other elements) then this is a finite cyclic group. Since this is not implemented directly in Sage as a group, *per se*, it is a bit harder to construct things like subgroups, but it is an excellent exercise to try. It is a nice example since the complex numbers are a concrete and familiar construction. Here are a few sample calculations to provide you with some exploratory tools. See the notes following the computations.

```
G = CyclotomicField(14)
w = G.gen(0); w
```

```
zeta14
```

```
wc = CDF(w)
wc.abs()
```

```
1.0
```

```
wc.arg()/N(2*pi/14)
```

```
1.0
```

```
b = w^2
b.multiplicative_order()
```

```
7
```

```
bc = CDF(b); bc
```

```
0.62348980185... + 0.781831482468...*I
```

```
bc.abs()
```

```
1.0
```

```
bc.arg()/N(2*pi/14)
```

```
2.0
```

```
sg = [b^i for i in range(7)]; sg
```

```
[1, zeta14^2, zeta14^4,
zeta14^5 - zeta14^4 + zeta14^3 - zeta14^2 + zeta14 - 1,
-zeta14, -zeta14^3, -zeta14^5]
```

```
c = sg[3]; d = sg[5]
c*d
```

```
zeta14^2
```

```
c = sg[3]; d = sg[6]
c*d in sg
```

```
True
```

```
c*d == sg[2]
```

```
True
```

```
sg[5]*sg[6] == sg[4]
```

```
True
```

```
G.multiplication_table(elements=sg)
```

```
*  a b c d e f g
 +--------------
a| a b c d e f g
b| b c d e f g a
c| c d e f g a b
d| d e f g a b c
e| e f g a b c d
f| f g a b c d e
g| g a b c d e f
```

Notes:

1. zeta14 is the name of the generator used for the cyclotomic field, it is a primitive root of unity (a 14th root of unity in this case). We have captured it as w.

2. The syntax CDF(w) will convert the complex number w into the more familiar form with real and imaginary parts.

3. The method .abs() will return the modulus of a complex number, r as described in the text. For elements of $\mathbb{C}^*$ this should always equal 1.

4. The method `.arg()` will return the argument of a complex number, θ as described in the text. Every element of the cyclic group in this example should have an argument that is an integer multiple of $\frac{2\pi}{14}$. The `N()` syntax converts the symbolic value of `pi` to a numerical approximation.

5. `sg` is a list of elements that form a cyclic subgroup of order 7, composed of the first 7 powers of `b = w^2`. So, for example, the last comparison multiplies the fifth power of `b` with the sixth power of `b`, which would be the eleventh power of `b`. But since `b` has order 7, this reduces to the fourth power.

6. If you know a subset of an infinite group forms a subgroup, then you can produce its Cayley table by specifying the list of elements you want to use. Here we ask for a multiplication table, since that is the relevant operation.

4.7   Sage Exercises

This group of exercises is about the group of units mod n, $U(n)$, which is sometimes cyclic, sometimes not. There are some commands in Sage that will answer some of these questions very quickly, but instead of using those now, just use the basic techniques described. The idea here is to just work with elements, and lists of elements, to discern the subgroup structure of these groups.

Sage worksheets have extensive capabilities for making new cells with carefully formatted text, include support for LATEX syntax to express mathematics. So when a question asks for explanation or commentary, make a new cell and communicate clearly with your audience. Continue this practice in subsequent exercise sets.**1**.   Execute the statement `R = Integers(40)` to create the set `[0,1,2,...,39]` This is a group under addition mod 40, which we will ignore. Instead we are interested in the subset of elements which have an inverse under *multiplication* mod 40. Determine how big this subgroup is by executing the command `R.unit_group_order()`, and then obtain a list of these elements with `R.list_of_elements_of_multiplicative_group()`.

2.   You can create elements of this group by coercing regular integers into U, such as with the statement `a = U(7)`. (Don't confuse this with our mathematical notation $U(40)$.) This will tell Sage that you want to view 7 as an element of U, subject to the corresponding operations. Determine the elements of the cyclic subgroup of U generated by 7 with a list comprehension as follows:

```
R = Integers(40)
a = R(7)
[a^i for i in srange(16)]
```

What is the order of 7 in $U(40)$?

3.   The group $U(49)$ is cyclic. Using only the Sage commands described previously, use Sage to find a generator for this group. Now using *only* theorems about the structure of cyclic groups, describe each of the subgroups of $U(49)$ by specifying its order and by giving an explicit generator. Do not repeat any of the subgroups — in other words, present each subgroup *exactly* once. You can use Sage to check your work on the subgroups, but your answer about the subgroups should rely only on theorems and be a nicely written paragraph with a table, etc.

4.   The group $U(35)$ is not cyclic. Again, using only the Sage commands described previously, use computations to provide irrefutable evidence of this. How many of the 16 different subgroups of $U(35)$ can you list?

5.   Again, using only the Sage commands described previously, explore the structure of $U(n)$ for various values of n and see if you can formulate an interesting conjecture about some basic property of this group. (Yes, this is a *very* open-ended question, but this is ultimately the real power of exploring mathematics with Sage.)

Permutation Groups

Permutation groups are central to the study of geometric symmetries and to Galois theory, the study of finding solutions of polynomial equations. They also provide abundant examples of nonabelian groups.

Let us recall for a moment the symmetries of the equilateral triangle $\triangle ABC$ from Chapter 3. The symmetries actually consist of permutations of the three vertices, where a **permutation** of the set $S = \{A, B, C\}$ is a one-to-one and onto map $\pi : S \to S$. The three vertices have the following six permutations.

$$\begin{pmatrix} A & B & C \\ A & B & C \end{pmatrix} \qquad \begin{pmatrix} A & B & C \\ C & A & B \end{pmatrix} \qquad \begin{pmatrix} A & B & C \\ B & C & A \end{pmatrix}$$

$$\begin{pmatrix} A & B & C \\ A & C & B \end{pmatrix} \qquad \begin{pmatrix} A & B & C \\ C & B & A \end{pmatrix} \qquad \begin{pmatrix} A & B & C \\ B & A & C \end{pmatrix}$$

We have used the array

$$\begin{pmatrix} A & B & C \\ B & C & A \end{pmatrix}$$

to denote the permutation that sends A to B, B to C, and C to A. That is,

$$A \mapsto B$$
$$B \mapsto C$$
$$C \mapsto A.$$

The symmetries of a triangle form a group. In this chapter we will study groups of this type.

5.1 Definitions and Notation

In general, the permutations of a set X form a group S_X. If X is a finite set, we can assume $X = \{1, 2, \ldots, n\}$. In this case we write S_n instead of S_X. The following theorem says that S_n is a group. We call this group the **symmetric group** on n letters.

Theorem 5.1. *The symmetric group on n letters, S_n, is a group with $n!$ elements, where the binary operation is the composition of maps.*

PROOF. The identity of S_n is just the identity map that sends 1 to 1, 2 to 2, ..., n to n. If $f : S_n \to S_n$ is a permutation, then f^{-1} exists, since f is one-to-one and onto; hence, every permutation has an inverse. Composition of maps is associative, which makes the group operation associative. We leave the proof that $|S_n| = n!$ as an exercise. $\qquad \square$

A subgroup of S_n is called a **permutation group**.

Example 5.2. Consider the subgroup G of S_5 consisting of the identity permutation id and the permutations

$$\sigma = \begin{pmatrix} 1 & 2 & 3 & 4 & 5 \\ 1 & 2 & 3 & 5 & 4 \end{pmatrix}$$

$$\tau = \begin{pmatrix} 1 & 2 & 3 & 4 & 5 \\ 3 & 2 & 1 & 4 & 5 \end{pmatrix}$$

$$\mu = \begin{pmatrix} 1 & 2 & 3 & 4 & 5 \\ 3 & 2 & 1 & 5 & 4 \end{pmatrix}.$$

The following table tells us how to multiply elements in the permutation group G.

$\circ$	id	σ	τ	μ
id	id	σ	τ	μ
σ	σ	id	μ	τ
τ	τ	μ	id	σ
μ	μ	τ	σ	id

Remark 5.3. Though it is natural to multiply elements in a group from left to right, functions are composed from right to left. Let σ and τ be permutations on a set X. To compose σ and τ as functions, we calculate $(\sigma \circ \tau)(x) = \sigma(\tau(x))$. That is, we do τ first, then σ. There are several ways to approach this inconsistency. *We will adopt the convention of multiplying permutations right to left. To compute $\sigma\tau$, do τ first and then σ.* That is, by $\sigma\tau(x)$ we mean $\sigma(\tau(x))$. (Another way of solving this problem would be to write functions on the right; that is, instead of writing $\sigma(x)$, we could write $(x)\sigma$. We could also multiply permutations left to right to agree with the usual way of multiplying elements in a group. Certainly all of these methods have been used.

Example 5.4. Permutation multiplication is not usually commutative. Let

$$\sigma = \begin{pmatrix} 1 & 2 & 3 & 4 \\ 4 & 1 & 2 & 3 \end{pmatrix}$$

$$\tau = \begin{pmatrix} 1 & 2 & 3 & 4 \\ 2 & 1 & 4 & 3 \end{pmatrix}.$$

Then

$$\sigma\tau = \begin{pmatrix} 1 & 2 & 3 & 4 \\ 1 & 4 & 3 & 2 \end{pmatrix},$$

but

$$\tau\sigma = \begin{pmatrix} 1 & 2 & 3 & 4 \\ 3 & 2 & 1 & 4 \end{pmatrix}.$$

Cycle Notation

The notation that we have used to represent permutations up to this point is cumbersome, to say the least. To work effectively with permutation groups, we need a more streamlined method of writing down and manipulating permutations.

A permutation $\sigma \in S_X$ is a **cycle of length** k if there exist elements $a_1, a_2, \ldots, a_k \in X$ such that

$$\sigma(a_1) = a_2$$

$$\sigma(a_2) = a_3$$

$$\vdots$$

$$\sigma(a_k) = a_1$$

and $\sigma(x) = x$ for all other elements $x \in X$. We will write $(a_1, a_2, \ldots, a_k)$ to denote the cycle σ. Cycles are the building blocks of all permutations.

Example 5.5. The permutation

$$\sigma = \begin{pmatrix} 1 & 2 & 3 & 4 & 5 & 6 & 7 \\ 6 & 3 & 5 & 1 & 4 & 2 & 7 \end{pmatrix} = (162354)$$

is a cycle of length 6, whereas

$$\tau = \begin{pmatrix} 1 & 2 & 3 & 4 & 5 & 6 \\ 1 & 4 & 2 & 3 & 5 & 6 \end{pmatrix} = (243)$$

is a cycle of length 3.

Not every permutation is a cycle. Consider the permutation

$$\begin{pmatrix} 1 & 2 & 3 & 4 & 5 & 6 \\ 2 & 4 & 1 & 3 & 6 & 5 \end{pmatrix} = (1243)(56).$$

This permutation actually contains a cycle of length 2 and a cycle of length 4.

Example 5.6. It is very easy to compute products of cycles. Suppose that

$$\sigma = (1352) \quad \text{and} \quad \tau = (256).$$

If we think of σ as

$$1 \mapsto 3, \qquad 3 \mapsto 5, \qquad 5 \mapsto 2, \qquad 2 \mapsto 1,$$

and τ as

$$2 \mapsto 5, \qquad 5 \mapsto 6, \qquad 6 \mapsto 2,$$

then for $\sigma\tau$ remembering that we apply τ first and then σ, it must be the case that

$$1 \mapsto 3, \qquad 3 \mapsto 5, \qquad 5 \mapsto 6, \qquad 6 \mapsto 2 \mapsto 1,$$

or $\sigma\tau = (1356)$. If $\mu = (1634)$, then $\sigma\mu = (1652)(34)$.

Two cycles in S_X, $\sigma = (a_1, a_2, \ldots, a_k)$ and $\tau = (b_1, b_2, \ldots, b_l)$, are ***disjoint*** if $a_i \neq b_j$ for all i and j.

Example 5.7. The cycles (135) and (27) are disjoint; however, the cycles (135) and (347) are not. Calculating their products, we find that

$$(135)(27) = (135)(27)$$
$$(135)(347) = (13475).$$

The product of two cycles that are not disjoint may reduce to something less complicated; the product of disjoint cycles cannot be simplified.

Proposition 5.8. *Let σ and τ be two disjoint cycles in S_X. Then $\sigma\tau = \tau\sigma$.*

PROOF.   Let $\sigma = (a_1, a_2, \ldots, a_k)$ and $\tau = (b_1, b_2, \ldots, b_l)$. We must show that $\sigma\tau(x) = \tau\sigma(x)$ for all $x \in X$. If x is neither in $\{a_1, a_2, \ldots, a_k\}$ nor $\{b_1, b_2, \ldots, b_l\}$, then both σ and τ fix x. That is, $\sigma(x) = x$ and $\tau(x) = x$. Hence,

$$\sigma\tau(x) = \sigma(\tau(x)) = \sigma(x) = x = \tau(x) = \tau(\sigma(x)) = \tau\sigma(x).$$

Do not forget that we are multiplying permutations right to left, which is the opposite of the order in which we usually multiply group elements. Now suppose that $x \in \{a_1, a_2, \ldots, a_k\}$. Then $\sigma(a_i) = a_{(i \bmod k)+1}$; that is,

$$a_1 \mapsto a_2$$
$$a_2 \mapsto a_3$$
$$\vdots$$
$$a_{k-1} \mapsto a_k$$
$$a_k \mapsto a_1.$$

However, $\tau(a_i) = a_i$ since σ and τ are disjoint. Therefore,

$$\begin{aligned}
\sigma\tau(a_i) &= \sigma(\tau(a_i)) \\
&= \sigma(a_i) \\
&= a_{(i \bmod k)+1} \\
&= \tau(a_{(i \bmod k)+1}) \\
&= \tau(\sigma(a_i)) \\
&= \tau\sigma(a_i).
\end{aligned}$$

Similarly, if $x \in \{b_1, b_2, \ldots, b_l\}$, then σ and τ also commute. $\square$

Theorem 5.9. *Every permutation in S_n can be written as the product of disjoint cycles.*

PROOF.   We can assume that $X = \{1, 2, \ldots, n\}$. If $\sigma \in S_n$ and we define X_1 to be $\{\sigma(1), \sigma^2(1), \ldots\}$, then the set X_1 is finite since X is finite. Now let i be the first integer in X that is not in X_1 and define X_2 by $\{\sigma(i), \sigma^2(i), \ldots\}$. Again, X_2 is a finite set. Continuing in this manner, we can define finite disjoint sets $X_3, X_4, \ldots$. Since X is a finite set, we are guaranteed that this process will end and there will be only a finite number of these sets, say r. If σ_i is the cycle defined by

$$\sigma_i(x) = \begin{cases} \sigma(x) & x \in X_i \\ x & x \notin X_i, \end{cases}$$

then $\sigma = \sigma_1 \sigma_2 \cdots \sigma_r$. Since the sets $X_1, X_2, \ldots, X_r$ are disjoint, the cycles $\sigma_1, \sigma_2, \ldots, \sigma_r$ must also be disjoint. $\square$

Example 5.10. Let

$$\sigma = \begin{pmatrix} 1 & 2 & 3 & 4 & 5 & 6 \\ 6 & 4 & 3 & 1 & 5 & 2 \end{pmatrix}$$
$$\tau = \begin{pmatrix} 1 & 2 & 3 & 4 & 5 & 6 \\ 3 & 2 & 1 & 5 & 6 & 4 \end{pmatrix}.$$

Using cycle notation, we can write

$$\sigma = (1624)$$

$$\tau = (13)(456)$$
$$\sigma\tau = (136)(245)$$
$$\tau\sigma = (143)(256).$$

Remark 5.11. From this point forward we will find it convenient to use cycle notation to represent permutations. When using cycle notation, we often denote the identity permutation by (1).

Transpositions

The simplest permutation is a cycle of length 2. Such cycles are called ***transpositions***. Since

$$(a_1, a_2, \ldots, a_n) = (a_1 a_n)(a_1 a_{n-1}) \cdots (a_1 a_3)(a_1 a_2),$$

any cycle can be written as the product of transpositions, leading to the following proposition.

Proposition 5.12. *Any permutation of a finite set containing at least two elements can be written as the product of transpositions.*

Example 5.13. Consider the permutation

$$(16)(253) = (16)(23)(25) = (16)(45)(23)(45)(25).$$

As we can see, there is no unique way to represent permutation as the product of transpositions. For instance, we can write the identity permutation as $(12)(12)$, as $(13)(24)(13)(24)$, and in many other ways. However, as it turns out, no permutation can be written as the product of both an even number of transpositions and an odd number of transpositions. For instance, we could represent the permutation (16) by

$$(23)(16)(23)$$

or by

$$(35)(16)(13)(16)(13)(35)(56),$$

but (16) will always be the product of an odd number of transpositions.

Lemma 5.14. *If the identity is written as the product of r transpositions,*

$$\mathrm{id} = \tau_1 \tau_2 \cdots \tau_r,$$

then r is an even number.

PROOF. We will employ induction on r. A transposition cannot be the identity; hence, $r > 1$. If $r = 2$, then we are done. Suppose that $r > 2$. In this case the product of the last two transpositions, $\tau_{r-1}\tau_r$, must be one of the following cases:

$$(ab)(ab) = \mathrm{id}$$
$$(bc)(ab) = (ac)(bc)$$
$$(cd)(ab) = (ab)(cd)$$
$$(ac)(ab) = (ab)(bc),$$

where a, b, c, and d are distinct.

The first equation simply says that a transposition is its own inverse. If this case occurs, delete $\tau_{r-1}\tau_r$ from the product to obtain

$$\mathrm{id} = \tau_1 \tau_2 \cdots \tau_{r-3} \tau_{r-2}.$$

By induction $r - 2$ is even; hence, r must be even.

In each of the other three cases, we can replace $\tau_{r-1}\tau_r$ with the right-hand side of the corresponding equation to obtain a new product of r transpositions for the identity. In this new product the last occurrence of a will be in the next-to-the-last transposition. We can continue this process with $\tau_{r-2}\tau_{r-1}$ to obtain either a product of $r - 2$ transpositions or a new product of r transpositions where the last occurrence of a is in τ_{r-2}. If the identity is the product of $r - 2$ transpositions, then again we are done, by our induction hypothesis; otherwise, we will repeat the procedure with $\tau_{r-3}\tau_{r-2}$.

At some point either we will have two adjacent, identical transpositions canceling each other out or a will be shuffled so that it will appear only in the first transposition. However, the latter case cannot occur, because the identity would not fix a in this instance. Therefore, the identity permutation must be the product of $r - 2$ transpositions and, again by our induction hypothesis, we are done. $\square$

Theorem 5.15. *If a permutation σ can be expressed as the product of an even number of transpositions, then any other product of transpositions equaling σ must also contain an even number of transpositions. Similarly, if σ can be expressed as the product of an odd number of transpositions, then any other product of transpositions equaling σ must also contain an odd number of transpositions.*

PROOF.   Suppose that

$$\sigma = \sigma_1\sigma_2\cdots\sigma_m = \tau_1\tau_2\cdots\tau_n,$$

where m is even. We must show that n is also an even number. The inverse of σ is $\sigma_m\cdots\sigma_1$. Since

$$\text{id} = \sigma\sigma_m\cdots\sigma_1 = \tau_1\cdots\tau_n\sigma_m\cdots\sigma_1,$$

n must be even by Lemma 5.14. The proof for the case in which σ can be expressed as an odd number of transpositions is left as an exercise. $\square$

In light of Theorem 5.15, we define a permutation to be **even** if it can be expressed as an even number of transpositions and **odd** if it can be expressed as an odd number of transpositions.

The Alternating Groups

One of the most important subgroups of S_n is the set of all even permutations, A_n. The group A_n is called the **alternating group on n letters**.

Theorem 5.16. *The set A_n is a subgroup of S_n.*

PROOF.   Since the product of two even permutations must also be an even permutation, A_n is closed. The identity is an even permutation and therefore is in A_n. If σ is an even permutation, then

$$\sigma = \sigma_1\sigma_2\cdots\sigma_r,$$

where σ_i is a transposition and r is even. Since the inverse of any transposition is itself,

$$\sigma^{-1} = \sigma_r\sigma_{r-1}\cdots\sigma_1$$

is also in A_n. $\square$

Proposition 5.17. *The number of even permutations in S_n, $n \geq 2$, is equal to the number of odd permutations; hence, the order of A_n is $n!/2$.*

PROOF. Let A_n be the set of even permutations in S_n and B_n be the set of odd permutations. If we can show that there is a bijection between these sets, they must contain the same number of elements. Fix a transposition σ in S_n. Since $n \geq 2$, such a σ exists. Define

$$\lambda_\sigma : A_n \to B_n$$

by

$$\lambda_\sigma(\tau) = \sigma\tau.$$

Suppose that $\lambda_\sigma(\tau) = \lambda_\sigma(\mu)$. Then $\sigma\tau = \sigma\mu$ and so

$$\tau = \sigma^{-1}\sigma\tau = \sigma^{-1}\sigma\mu = \mu.$$

Therefore, λ_σ is one-to-one. We will leave the proof that λ_σ is surjective to the reader. □

Example 5.18. The group A_4 is the subgroup of S_4 consisting of even permutations. There are twelve elements in A_4:

(1)	(12)(34)	(13)(24)	(14)(23)
(123)	(132)	(124)	(142)
(134)	(143)	(234)	(243).

One of the end-of-chapter exercises will be to write down all the subgroups of A_4. You will find that there is no subgroup of order 6. Does this surprise you?

--- ***Historical Note*** ---

Lagrange first thought of permutations as functions from a set to itself, but it was Cauchy who developed the basic theorems and notation for permutations. He was the first to use cycle notation. Augustin-Louis Cauchy (1789–1857) was born in Paris at the height of the French Revolution. His family soon left Paris for the village of Arcueil to escape the Reign of Terror. One of the family's neighbors there was Pierre-Simon Laplace (1749–1827), who encouraged him to seek a career in mathematics. Cauchy began his career as a mathematician by solving a problem in geometry given to him by Lagrange. Cauchy wrote over 800 papers on such diverse topics as differential equations, finite groups, applied mathematics, and complex analysis. He was one of the mathematicians responsible for making calculus rigorous. Perhaps more theorems and concepts in mathematics have the name Cauchy attached to them than that of any other mathematician.

5.2 Dihedral Groups

Another special type of permutation group is the dihedral group. Recall the symmetry group of an equilateral triangle in Chapter 3. Such groups consist of the rigid motions of a regular n-sided polygon or n-gon. For $n = 3, 4, \ldots$, we define the ***nth dihedral group*** to be the group of rigid motions of a regular n-gon. We will denote this group by D_n. We can number the vertices of a regular n-gon by $1, 2, \ldots, n$ (Figure 5.19). Notice that there are exactly n choices to replace the first vertex. If we replace the first vertex by k, then the second vertex must be replaced either by vertex $k+1$ or by vertex $k-1$; hence, there are $2n$ possible rigid motions of the n-gon. We summarize these results in the following theorem.
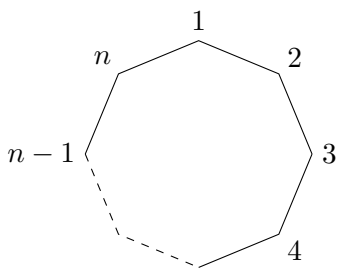
Figure 5.19: A regular n-gon

Theorem 5.20. *The dihedral group, D_n, is a subgroup of S_n of order $2n$.*
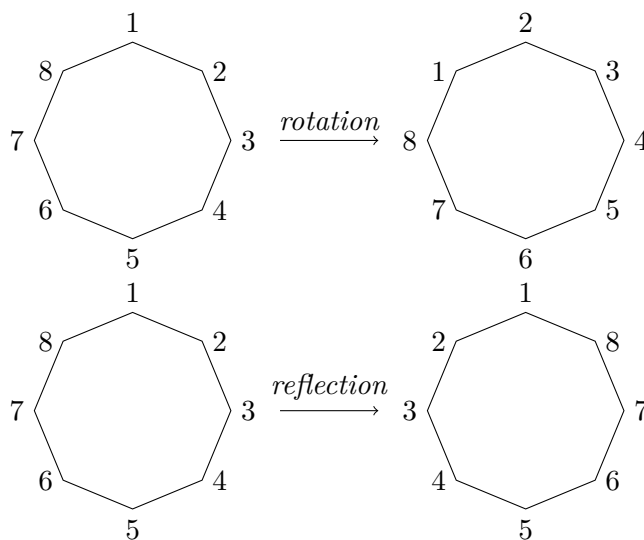
Figure 5.21: Rotations and reflections of a regular n-gon
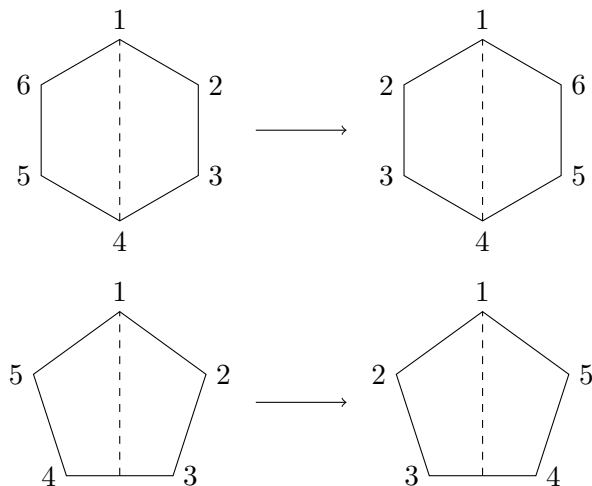
Figure 5.22: Types of reflections of a regular n-gon

Theorem 5.23. *The group D_n, $n \geq 3$, consists of all products of the two elements r and s, satisfying the relations*

$$r^n = 1$$

$$s^2 = 1$$
$$srs = r^{-1}.$$

PROOF. The possible motions of a regular n-gon are either reflections or rotations (Figure 5.21). There are exactly n possible rotations:

$$\text{id}, \frac{360°}{n}, 2 \cdot \frac{360°}{n}, \ldots, (n-1) \cdot \frac{360°}{n}.$$

We will denote the rotation $360°/n$ by r. The rotation r generates all of the other rotations. That is,

$$r^k = k \cdot \frac{360°}{n}.$$

Label the n reflections $s_1, s_2, \ldots, s_n$, where s_k is the reflection that leaves vertex k fixed. There are two cases of reflections, depending on whether n is even or odd. If there are an even number of vertices, then two vertices are left fixed by a reflection, and $s_1 = s_{n/2+1}, s_2 = s_{n/2+2}, \ldots, s_{n/2} = s_n$. If there are an odd number of vertices, then only a single vertex is left fixed by a reflection and $s_1, s_2, \ldots, s_n$ are distinct (Figure 5.22). In either case, the order of each s_k is two. Let $s = s_1$. Then $s^2 = 1$ and $r^n = 1$. Since any rigid motion t of the n-gon replaces the first vertex by the vertex k, the second vertex must be replaced by either $k+1$ or by $k-1$. If the second vertex is replaced by $k+1$, then $t = r^k$. If the second vertex is replaced by $k-1$, then $t = sr^k$. Hence, r and s generate D_n. That is, D_n consists of all finite products of r and s,

$$D_n = \{1, r, r^2, \ldots, r^{n-1}, s, sr, sr^2, \ldots, sr^{n-1}\}.$$

We will leave the proof that $srs = r^{-1}$ as an exercise. □

Example 5.24. The group of rigid motions of a square, D_4, consists of eight elements. With the vertices numbered 1, 2, 3, 4 (Figure 5.25), the rotations are

$$r = (1234)$$
$$r^2 = (13)(24)$$
$$r^3 = (1432)$$
$$r^4 = (1)$$

and the reflections are

$$s_1 = (24)$$
$$s_2 = (13).$$

The order of D_4 is 8. The remaining two elements are

$$rs_1 = (12)(34)$$
$$r^3 s_1 = (14)(23).$$
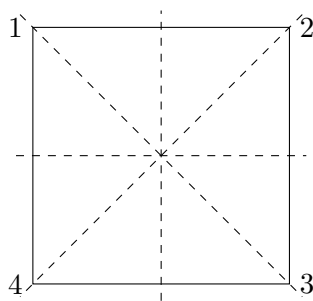
Figure 5.25: The group D_4

The Motion Group of a Cube

We can investigate the groups of rigid motions of geometric objects other than a regular n-sided polygon to obtain interesting examples of permutation groups. Let us consider the group of rigid motions of a cube. One of the first questions that we can ask about this group is "what is its order?" A cube has 6 sides. If a particular side is facing upward, then there are four possible rotations of the cube that will preserve the upward-facing side. Hence, the order of the group is $6 \cdot 4 = 24$. We have just proved the following proposition.
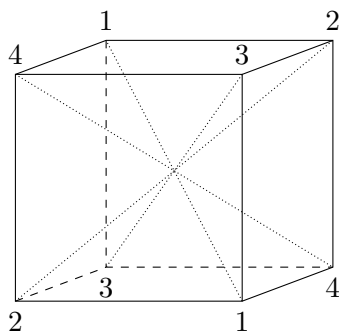
Figure 5.26: The motion group of a cube

Proposition 5.27. *The group of rigid motions of a cube contains* 24 *elements.*

Theorem 5.28. *The group of rigid motions of a cube is S_4.*

PROOF.   From Proposition 5.27, we already know that the motion group of the cube has 24 elements, the same number of elements as there are in S_4. There are exactly four diagonals in the cube. If we label these diagonals 1, 2, 3, and 4, we must show that the motion group of the cube will give us any permutation of the diagonals (Figure 5.26). If we can obtain all of these permutations, then S_4 and the group of rigid motions of the cube must be the same. To obtain a transposition we can rotate the cube 180° about the axis joining the midpoints of opposite edges (Figure 5.29). There are six such axes, giving all transpositions in S_4. Since every element in S_4 is the product of a finite number of transpositions, the motion group of a cube must be S_4.   □
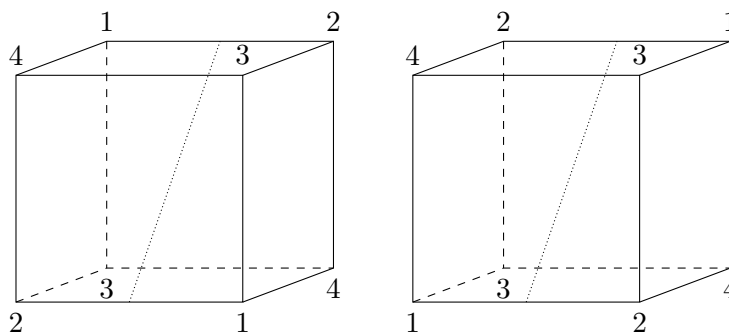
Figure 5.29: Transpositions in the motion group of a cube

5.3 Exercises

1. Write the following permutations in cycle notation.

(a)
$$\begin{pmatrix} 1 & 2 & 3 & 4 & 5 \\ 2 & 4 & 1 & 5 & 3 \end{pmatrix}$$

(c)
$$\begin{pmatrix} 1 & 2 & 3 & 4 & 5 \\ 3 & 5 & 1 & 4 & 2 \end{pmatrix}$$

(b)
$$\begin{pmatrix} 1 & 2 & 3 & 4 & 5 \\ 4 & 2 & 5 & 1 & 3 \end{pmatrix}$$

(d)
$$\begin{pmatrix} 1 & 2 & 3 & 4 & 5 \\ 1 & 4 & 3 & 2 & 5 \end{pmatrix}$$

2. Compute each of the following.

(a) $(1345)(234)$

(b) $(12)(1253)$

(c) $(143)(23)(24)$

(d) $(1423)(34)(56)(1324)$

(e) $(1254)(13)(25)$

(f) $(1254)(13)(25)^2$

(g) $(1254)^{-1}(123)(45)(1254)$

(h) $(1254)^2(123)(45)$

(i) $(123)(45)(1254)^{-2}$

(j) $(1254)^{100}$

(k) $|(1254)|$

(l) $|(1254)^2|$

(m) $(12)^{-1}$

(n) $(12537)^{-1}$

(o) $[(12)(34)(12)(47)]^{-1}$

(p) $[(1235)(467)]^{-1}$

3. Express the following permutations as products of transpositions and identify them as even or odd.

(a) (14356)

(b) $(156)(234)$

(c) $(1426)(142)$

(d) $(17254)(1423)(154632)$

(e) (142637)

4. Find $(a_1, a_2, \ldots, a_n)^{-1}$.

5. List all of the subgroups of S_4. Find each of the following sets:

(a) $\{\sigma \in S_4 : \sigma(1) = 3\}$

(b) $\{\sigma \in S_4 : \sigma(2) = 2\}$

(c) $\{\sigma \in S_4 : \sigma(1) = 3 \text{ and } \sigma(2) = 2\}$.

Are any of these sets subgroups of S_4?

6. Find all of the subgroups in A_4. What is the order of each subgroup?

7. Find all possible orders of elements in S_7 and A_7.

8. Show that A_{10} contains an element of order 15.

9. Does A_8 contain an element of order 26?

10. Find an element of largest order in S_n for $n = 3, \ldots, 10$.

11. What are the possible cycle structures of elements of A_5? What about A_6?

12. Let $\sigma \in S_n$ have order n. Show that for all integers i and j, $\sigma^i = \sigma^j$ if and only if $i \equiv j \pmod{n}$.

13. Let $\sigma = \sigma_1 \cdots \sigma_m \in S_n$ be the product of disjoint cycles. Prove that the order of σ is the least common multiple of the lengths of the cycles $\sigma_1, \ldots, \sigma_m$.

14. Using cycle notation, list the elements in D_5. What are r and s? Write every element as a product of r and s.

15. If the diagonals of a cube are labeled as Figure 5.26, to which motion of the cube does the permutation $(12)(34)$ correspond? What about the other permutations of the diagonals?

16. Find the group of rigid motions of a tetrahedron. Show that this is the same group as A_4.

17. Prove that S_n is nonabelian for $n \geq 3$.

18. Show that A_n is nonabelian for $n \geq 4$.

19. Prove that D_n is nonabelian for $n \geq 3$.

20. Let $\sigma \in S_n$ be a cycle. Prove that σ can be written as the product of at most $n - 1$ transpositions.

21. Let $\sigma \in S_n$. If σ is not a cycle, prove that σ can be written as the product of at most $n - 2$ transpositions.

22. If σ can be expressed as an odd number of transpositions, show that any other product of transpositions equaling σ must also be odd.

23. If σ is a cycle of odd length, prove that σ^2 is also a cycle.

24. Show that a 3-cycle is an even permutation.

25. Prove that in A_n with $n \geq 3$, any permutation is a product of cycles of length 3.

26. Prove that any element in S_n can be written as a finite product of the following permutations.

(a) $(12), (13), \ldots, (1n)$

(b) $(12), (23), \ldots, (n - 1, n)$

(c) $(12), (12 \ldots n)$

27. Let G be a group and define a map $\lambda_g : G \to G$ by $\lambda_g(a) = ga$. Prove that λ_g is a permutation of G.

28. Prove that there exist $n!$ permutations of a set containing n elements.

29. Recall that the **center** of a group G is

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

Find the center of D_8. What about the center of D_{10}? What is the center of D_n?

30. Let $\tau = (a_1, a_2, \ldots, a_k)$ be a cycle of length k.

(a) Prove that if σ is any permutation, then

$$\sigma\tau\sigma^{-1} = (\sigma(a_1), \sigma(a_2), \ldots, \sigma(a_k))$$

is a cycle of length k.

(b) Let μ be a cycle of length k. Prove that there is a permutation σ such that $\sigma\tau\sigma^{-1} = \mu$.

31. For α and β in S_n, define $\alpha \sim \beta$ if there exists an $\sigma \in S_n$ such that $\sigma\alpha\sigma^{-1} = \beta$. Show that $\sim$ is an equivalence relation on S_n.

32. Let $\sigma \in S_X$. If $\sigma^n(x) = y$, we will say that $x \sim y$.

(a) Show that $\sim$ is an equivalence relation on X.

(b) If $\sigma \in A_n$ and $\tau \in S_n$, show that $\tau^{-1}\sigma\tau \in A_n$.

(c) Define the **orbit** of $x \in X$ under $\sigma \in S_X$ to be the set

$$\mathcal{O}_{x,\sigma} = \{y : x \sim y\}.$$

Compute the orbits of each of the following elements in S_5:

$$\alpha = (1254)$$
$$\beta = (123)(45)$$
$$\gamma = (13)(25).$$

(d) If $\mathcal{O}_{x,\sigma} \cap \mathcal{O}_{y,\sigma} \neq \emptyset$, prove that $\mathcal{O}_{x,\sigma} = \mathcal{O}_{y,\sigma}$. The orbits under a permutation σ are the equivalence classes corresponding to the equivalence relation $\sim$.

(e) A subgroup H of S_X is **transitive** if for every $x, y \in X$, there exists a $\sigma \in H$ such that $\sigma(x) = y$. Prove that $\langle \sigma \rangle$ is transitive if and only if $\mathcal{O}_{x,\sigma} = X$ for some $x \in X$.

33. Let $\alpha \in S_n$ for $n \geq 3$. If $\alpha\beta = \beta\alpha$ for all $\beta \in S_n$, prove that α must be the identity permutation; hence, the center of S_n is the trivial subgroup.

34. If α is even, prove that α^{-1} is also even. Does a corresponding result hold if α is odd?

35. Show that $\alpha^{-1}\beta^{-1}\alpha\beta$ is even for $\alpha, \beta \in S_n$.

36. Let r and s be the elements in D_n described in Theorem 5.23

(a) Show that $srs = r^{-1}$.

(b) Show that $r^k s = sr^{-k}$ in D_n.

(c) Prove that the order of $r^k \in D_n$ is $n/\gcd(k, n)$.

5.4   Sage

A good portion of Sage's support for group theory is based on routines from GAP (Groups, Algorithms, and Programming) at www.gap-system.org, which is included in every copy of Sage. This is a mature open source package, dating back to 1986. (Forward reference here to GAP console, etc.)

As we have seen, groups can be described in many different ways, such as sets of matrices, sets of complex numbers, or sets of symbols subject to defining relations. A very concrete way to represent groups is via permutations (one-to-one and onto functions of the integers 1 through n), using function composition as the operation in the group, as described in this chapter. Sage has many routines designed to work with groups of this type and they are also a good way for those learning group theory to gain experience with the basic ideas of group theory. For both these reasons, we will concentrate on these types of groups.

Permutation Groups and Elements

The easiest way to work with permutation group elements in Sage is to write them in cycle notation. Since these are products of disjoint cycles (which commute), we do not need to concern ourselves with the actual order of the cycles. If we write `(1,3)(2,4)` we probably understand it to be a permutation (the topic of this chapter!) and we know that it could be an element of S_4, or perhaps a symmetric group on more symbols than just 4. Sage cannot get started that easily and needs a bit of context, so we coerce a string of characters written with cycle notation into a symmetric group to make group elements. Here are some examples and some sample computations. Remember that Sage and your text differ on how to interpret the order of composing two permutations in a product.

```
G = SymmetricGroup(5)
sigma = G("(1,3)(2,5,4)")
sigma*sigma
```

```
(2,4,5)
```

```
rho = G("(2,4)(1,5)")
rho^3
```

```
(1,5)(2,4)
```

If the next three examples seem confusing, or "backwards", then now would be an excellent time to review the Sage discussion about the order of permutation composition in the subsection Groups of symmetries.

```
sigma*rho
```

```
(1,3,5,2)
```

```
rho*sigma
```

```
(1,4,5,3)
```

```
rho^-1*sigma*rho
```

```
(1,2,4)(3,5)
```

There are alternate ways to create permutation group elements, which can be useful in some situations, but they are not quite as useful in everday use.

```
sigma1 = G("(1,3)(2,5,4)")
sigma1
```

(1,3)(2,5,4)

```
sigma2 = G([(1,3),(2,5,4)])
sigma2
```

(1,3)(2,5,4)

```
sigma3 = G([3,5,1,2,4])
sigma3
```

(1,3)(2,5,4)

```
sigma1 == sigma2
```

True

```
sigma2 == sigma3
```

True

```
sigma2.cycle_tuples()
```

[(1, 3), (2, 5, 4)]

```
[sigma3(x) for x in G.domain()]
```

[3, 5, 1, 2, 4]

The second version of σ is a list of "tuples", which requires a lot of commas and these must be enclosed in a list. (A tuple of length one must be written like (4,) to distinguish it from using parentheses for grouping, as in 5*(4).) The third version uses the "bottom-row" of the more cumbersome two-row notation introduced at the beginning of the chapter — it is an ordered list of the *output values* of the permutation when considered as a function.

So we then see that despite three different input procedures, all the versions of σ print the same way, and moreso they are actually equal to each other. (This is a subtle difference — what an object *is* in Sage versus how an object *displays* itself.)

We can be even more careful about the nature of our elements. Notice that once we get Sage started, it can promote the product $\tau\sigma$ into the larger permutation group. We can "promote" elements into larger permutation groups, but it is an error to try to shoe-horn an element into a too-small symmetric group.

```
H = SymmetricGroup(4)
sigma = H("(1,2,3,4)")
G = SymmetricGroup(6)
tau = G("(1,2,3,4,5,6)")
rho = tau * sigma
rho
```

(1,3)(2,4,5,6)

```
sigma.parent()
```

Symmetric group of order 4! as a permutation group

```
tau.parent()
```

Symmetric group of order 6! as a permutation group

```
rho.parent()
```

Symmetric group of order 6! as a permutation group

```
tau.parent() == rho.parent()
```

True

```
sigmaG = G(sigma)
sigmaG.parent()
```

Symmetric group of order 6! as a permutation group

It is an error to try to coerce a permutation with too many symbols into a permutation group employing too few symbols.

```
tauH = H(tau)
```

```
Traceback (most recent call last):
...
ValueError: Invalid permutation vector: (1,2,3,4,5,6)
```

Better than working with just elements of the symmetric group, we can create a variety of permutation groups in Sage. Here is a sampling for starters:

Sage Command	Description
SymmetricGroup(n)	Permutations on n symbols, $n!$ elements
DihedralGroup(n)	Symmetries of an n-gon, $2n$ elements.
CyclicPermutationGroup(n)	Rotations of an n-gon (no flips), n elements
AlternatingGroup(n)	Alternating group on n symbols, $n!/2$ elements
KleinFourGroup()	A non-cyclic group of order 4

Table 5.30: Some Sage permutation groups

You can also locate Sage permutation groups with the `groups` catalog. In the next cell place your cursor right after the final dot and hit the tab-key. You will get a list of methods you can use to create permutation groups. As always, place a question-mark after a method and hit the tab-key to get online documentation of a method.

```
groups.permutation.
```

Properties of Permutation Elements

Sometimes it is easier to grab an element out of a list of elements of a permutation group, and then it is already attached to a parent and there is no need for any coercion. In the following, `rotate` and `flip` are automatically elements of `G` because of the way we procured them.

```
D = DihedralGroup(5)
elements = D.list(); elements
```

```
[(), (1,5)(2,4), (1,2,3,4,5), (1,4)(2,3), (1,3,5,2,4), (2,5)(3,4),
 (1,3)(4,5), (1,5,4,3,2), (1,4,2,5,3), (1,2)(3,5)]
```

```
rotate = elements[2]
flip = elements[3]
flip*rotate == rotate* flip
```

```
False
```

So we see from this final statement that the group of symmetries of a pentagon is not abelian. But there is an easier way.

```
D = DihedralGroup(5)
D.is_abelian()
```

```
False
```

There are many more methods you can use for both permutation groups and their individual elements. Use the blank compute cell below to create a permutation group (any one you like) and an element of a permutation group (any one you like). Then use tab-completion to see all the methods available for an element, or for a group (name, period, tab-key). Some names you may recognize, some we will learn about in the coming chapters, some are highly-specialized research tools you can use when you write your Ph.D. thesis in group theory. For any of these methods, remember that you can type the name, followed by a question mark, to see documentation and examples. *Experiment and explore* — it is really hard to break anything.Here are some selected examples of various methods available.

```
A4 = AlternatingGroup(4)
A4.order()
```

```
12
```

```
A4.is_finite()
```

```
True
```

```
A4.is_abelian()
```

```
False
```

```
A4.is_cyclic()
```

```
False
```

```
sigma = A4("(1,2,4)")
sigma^-1
```

```
(1,4,2)
```

```
sigma.order()
```

```
3
```

A very useful method when studying the alternating group is the permutation group element method `.sign()`. It will return 1 if a permutation is even and `-1` if a permutation is odd.

```
G = SymmetricGroup(3)
sigma = G("(1,2)")
tau = G("(1,3)")
rho = sigma*tau
sigma.sign()
```

-1

```
rho.sign()
```

1

We can create subgroups by giving the main group a list of "generators." These elements serve to "generate" a subgroup — imagine multiplying these elements (and their inverses) together over and over, creating new elements that must also be in the subgroup and also become involved in new products, until you see no new elements. Now that definition ends with a horribly imprecise statement, but it should suffice for now. A better definition is that the subgroup generated by the elements is the smallest subgroup of the main group that contains all the generators — which is fine if you know what all the subgroups might be.

   With a single generator, the repeated products just become powers of the lone generator. The subgroup generated then is cyclic. With two (or more) generators, especially in a non-abelian group, the situation can be much, much more complicated. So let us begin with just a single generator. But do not forget to put it in a list anyway.

```
A4 = AlternatingGroup(4)
sigma = A4("(1,2,4)")
sg = A4.subgroup([sigma])
sg
```

```
Subgroup of (Alternating group of order 4!/2 as a permutation group)
generated by [(1,2,4)]
```

```
sg.order()
```

3

```
sg.list()
```

```
[(), (1,2,4), (1,4,2)]
```

```
sg.is_abelian()
```

True

```
sg.is_cyclic()
```

True

```
sg.is_subgroup(A4)
```

True

We can now redo the example from the very beginning of this chapter. We translate to elements to cycle notation, construct the subgroup from two generators (the subgroup is not cyclic), and since the subgroup is abelian, we do not have to view Sage's Cayley table as a diagonal reflection of the table in the example.

```
G = SymmetricGroup(5)
sigma = G("(4,5)")
tau = G("(1,3)")
H = G.subgroup([sigma, tau])
H.list()
```

```
[(), (1,3), (4,5), (1,3)(4,5)]
```

```
text_names = ['id', 'sigma', 'tau', 'mu']
H.cayley_table(names=text_names)
```

```
    *      id sigma   tau    mu
   +------------------------
   id|     id sigma   tau    mu
sigma| sigma    id    mu   tau
  tau|   tau    mu    id sigma
   mu|    mu   tau sigma    id
```

Motion Group of a Cube

We could mimic the example in the text and create elements of S_4 as permutations of the diagonals. A more obvious, but less insightful, construction is to view the 8 corners of the cube as the items being permuted. Then some obvious symmetries of the cube come from running an axis through the center of a side, through to the center of the opposite side, with quarter-turns or half-turns about these axes forming symmetries. With three such axes and four rotations per axis, we get 12 symmetries, except we have counted the identity permutation two extra times.

Label the four corners of the square top with 1 through 4, placing 1 in the left-front corner, and following around clockwise when viewed from above. Use 5 through 8 for the bottom square's corner, so that 5 is directly below 1, 6 below 2, etc. We will use quarter-turns, clockwise, around each axis, when viewed from above, the front, and the right.

```
G = SymmetricGroup(8)
above = G("(1,2,3,4)(5,6,7,8)")
front = G("(1,4,8,5)(2,3,7,6)")
right = G("(1,2,6,5)(3,7,8,4)")
cube = G.subgroup([above, front, right])
cube.order()
```

```
24
```

```
cube.list()
```

```
[(), (1,2,3,4)(5,6,7,8), (1,2,6,5)(3,7,8,4),
 (1,4,8,5)(2,3,7,6), (1,6,8)(2,7,4), (2,4,5)(3,8,6),
 (1,3,8)(2,7,5), (1,6)(2,5)(3,8)(4,7), (1,3,6)(4,7,5),
 (1,3)(2,4)(5,7)(6,8), (1,8)(2,7)(3,6)(4,5), (1,7)(2,3)(4,6)(5,8),
 (1,4)(2,8)(3,5)(6,7), (1,5,6,2)(3,4,8,7), (1,5,8,4)(2,6,7,3),
 (1,7)(2,6)(3,5)(4,8), (1,7)(2,8)(3,4)(5,6), (1,4,3,2)(5,8,7,6),
 (1,5)(2,8)(3,7)(4,6), (1,2)(3,5)(4,6)(7,8), (1,8,6)(2,4,7),
 (2,5,4)(3,6,8), (1,6,3)(4,5,7), (1,8,3)(2,5,7)]
```

Since we know from the discussion in the text that the symmetry group has 24 elements, we see that our three quarter-turns are sufficient to create every symmetry. This prompts several questions which you can find in Exercise 5.5.4.

5.5   Sage Exercises

These exercises are designed to help you become familiar with permutation groups in Sage.**1**. Create the full symmetric group S_{10} with the command G = SymmetricGroup(10).

2.   Create elements of G with the following (varying) syntax. Pay attention to commas, quotes, brackets, parentheses. The first two use a string (characters) as input, mimicking the way we write permuations (but with commas). The second two use a list of tuples.

- a = G("(5,7,2,9,3,1,8)")

- b = G("(1,3)(4,5)")

- c = G([(1,2),(3,4)])

- d = G([(1,3),(2,5,8),(4,6,7,9,10)])

(a) Compute a^3, bc, $ad^{-1}b$.

(b) Compute the orders of each of these four individual elements (a through d) using a single permutation group element method.

(c) Use the permutation group element method .sign() to determine if a, b, c, d are even or odd permutations.

(d) Create two cyclic subgroups of G with the commands:

- H = G.subgroup([a])
- K = G.subgroup([d])

List, and study, the elements of each subgroup. Without using Sage, list the order of each subgroup of K. Then use Sage to construct a subgroup of K with order 10.

(e) More complicated subgroups can be formed by using two or more generators. Construct a subgroup L of G with the command L = G.subgroup([b,c]). Compute the order of L and list all of the elements of L.

3.   Construct the group of symmetries of the tetrahedron (also the alternating group on 4 symbols, A_4) with the command A=AlternatingGroup(4). Using tools such as orders of elements, and generators of subgroups, see if you can find *all of* the subgroups of A_4 (each one exactly once). Do this without using the .subgroups() method to justify the correctness of your answer (though it might be a convenient way to check your work).

Provide a nice summary as your answer—not just piles of output. So use Sage as a tool, as needed, but basically your answer will be a concise paragraph and/or table. This is the one part of this assignment without clear, precise directions, so spend some time on this portion to get it right. Hint: no subgroup of A_4 requires more than two generators.

4.   The subsection Motion Group of a Cube describes the 24 symmetries of a cube as a subgroup of the symmetric group S_8 generated by three quarter-turns. Answer the following questions about this symmetry group.

(a) From the list of elements of the group, can you locate the ten rotations about axes? (Hint: the identity is easy, the other nine never send any symbol to itself.)

(b) Can you identify the six symmetries that are a transposition of diagonals? (Hint: `[g for g in cube if g.order() == 2]` is a good preliminary filter.)

(c) Verify that any two of the quarter-turns (`above`, `front`, `right`) are sufficient to generate the whole group. How do you know each pair generates the entire group?

(d) Can you express one of the diagonal transpositions as a product of quarter-turns? This can be a notoriously difficult problem, especially for software. It is known as the "word problem."

(e) Number the six faces of the cube with the numbers 1 through 6 (any way you like). Now consider the same three symmetries we used before (quarter-turns about face-to-face axes), but now view them as permutations of the six faces. In this way, we construct each symmetry as an element of S_6. Verify that the subgroup generated by these symmetries is the whole symmetry group of the cube. Again, rather than using three generators, try using just two.

5. Save your work, and then see if you can crash your Sage session by building the subgroup of S_{10} generated by the elements `b` and `d` of orders 2 and 30 from above. *Do not submit* the list of elements of `N` as part of your submitted worksheet.

```
N = G.subgroup([b,d])
N.list()
```

What is the order of N?

Cosets and Lagrange's Theorem

Lagrange's Theorem, one of the most important results in finite group theory, states that the order of a subgroup must divide the order of the group. This theorem provides a powerful tool for analyzing finite groups; it gives us an idea of exactly what type of subgroups we might expect a finite group to possess. Central to understanding Lagranges's Theorem is the notion of a coset.

6.1 Cosets

Let G be a group and H a subgroup of G. Define a ***left coset*** of H with ***representative*** $g \in G$ to be the set

$$gH = \{gh : h \in H\}.$$

Right cosets can be defined similarly by

$$Hg = \{hg : h \in H\}.$$

If left and right cosets coincide or if it is clear from the context to which type of coset that we are referring, we will use the word *coset* without specifying left or right.

Example 6.1. Let H be the subgroup of $\mathbb{Z}_6$ consisting of the elements 0 and 3. The cosets are

$$0 + H = 3 + H = \{0, 3\}$$
$$1 + H = 4 + H = \{1, 4\}$$
$$2 + H = 5 + H = \{2, 5\}.$$

We will always write the cosets of subgroups of $\mathbb{Z}$ and $\mathbb{Z}_n$ with the additive notation we have used for cosets here. In a commutative group, left and right cosets are always identical.

Example 6.2. Let H be the subgroup of S_3 defined by the permutations $\{(1), (123), (132)\}$. The left cosets of H are

$$(1)H = (123)H = (132)H = \{(1), (123), (132)\}$$
$$(12)H = (13)H = (23)H = \{(12), (13), (23)\}.$$

The right cosets of H are exactly the same as the left cosets:

$$H(1) = H(123) = H(132) = \{(1), (123), (132)\}$$
$$H(12) = H(13) = H(23) = \{(12), (13), (23)\}.$$

It is not always the case that a left coset is the same as a right coset. Let K be the subgroup of S_3 defined by the permutations $\{(1), (12)\}$. Then the left cosets of K are

$$
\begin{aligned}
(1)K &= (12)K = \{(1), (12)\} \\
(13)K &= (123)K = \{(13), (123)\} \\
(23)K &= (132)K = \{(23), (132)\};
\end{aligned}
$$

however, the right cosets of K are

$$
\begin{aligned}
K(1) &= K(12) = \{(1), (12)\} \\
K(13) &= K(132) = \{(13), (132)\} \\
K(23) &= K(123) = \{(23), (123)\}.
\end{aligned}
$$

The following lemma is quite useful when dealing with cosets. (We leave its proof as an exercise.)

Lemma 6.3. *Let H be a subgroup of a group G and suppose that $g_1, g_2 \in G$. The following conditions are equivalent.*

1. *$g_1 H = g_2 H$;*

2. *$H g_1^{-1} = H g_2^{-1}$;*

3. *$g_1 H \subset g_2 H$;*

4. *$g_2 \in g_1 H$;*

5. *$g_1^{-1} g_2 \in H$.*

In all of our examples the cosets of a subgroup H partition the larger group G. The following theorem proclaims that this will always be the case.

Theorem 6.4. *Let H be a subgroup of a group G. Then the left cosets of H in G partition G. That is, the group G is the disjoint union of the left cosets of H in G.*

PROOF.   Let $g_1 H$ and $g_2 H$ be two cosets of H in G. We must show that either $g_1 H \cap g_2 H = \emptyset$ or $g_1 H = g_2 H$. Suppose that $g_1 H \cap g_2 H \neq \emptyset$ and $a \in g_1 H \cap g_2 H$. Then by the definition of a left coset, $a = g_1 h_1 = g_2 h_2$ for some elements h_1 and h_2 in H. Hence, $g_1 = g_2 h_2 h_1^{-1}$ or $g_1 \in g_2 H$. By Lemma 6.3, $g_1 H = g_2 H$. $\square$

Remark 6.5. There is nothing special in this theorem about left cosets. Right cosets also partition G; the proof of this fact is exactly the same as the proof for left cosets except that all group multiplications are done on the opposite side of H.

Let G be a group and H be a subgroup of G. Define the **index** of H in G to be the number of left cosets of H in G. We will denote the index by $[G : H]$.

Example 6.6. Let $G = \mathbb{Z}_6$ and $H = \{0, 3\}$. Then $[G : H] = 3$.

Example 6.7. Suppose that $G = S_3$, $H = \{(1), (123), (132)\}$, and $K = \{(1), (12)\}$. Then $[G : H] = 2$ and $[G : K] = 3$.

Theorem 6.8. *Let H be a subgroup of a group G. The number of left cosets of H in G is the same as the number of right cosets of H in G.*

PROOF.   Let $\mathcal{L}_H$ and $\mathcal{R}_H$ denote the set of left and right cosets of H in G, respectively.   If we can define a bijective map $\phi : \mathcal{L}_H \to \mathcal{R}_H$, then the theorem will be proved. If $gH \in \mathcal{L}_H$, let $\phi(gH) = Hg^{-1}$. By Lemma 6.3, the map ϕ is well-defined; that is, if $g_1 H = g_2 H$, then $Hg_1^{-1} = Hg_2^{-1}$. To show that ϕ is one-to-one, suppose that

$$Hg_1^{-1} = \phi(g_1 H) = \phi(g_2 H) = Hg_2^{-1}.$$

Again by Lemma 6.3, $g_1 H = g_2 H$. The map ϕ is onto since $\phi(g^{-1}H) = Hg$.   □

6.2   Lagrange's Theorem

Proposition 6.9. *Let H be a subgroup of G with $g \in G$ and define a map $\phi : H \to gH$ by $\phi(h) = gh$. The map ϕ is bijective; hence, the number of elements in H is the same as the number of elements in gH.*

PROOF.   We first show that the map ϕ is one-to-one. Suppose that $\phi(h_1) = \phi(h_2)$ for elements $h_1, h_2 \in H$. We must show that $h_1 = h_2$, but $\phi(h_1) = gh_1$ and $\phi(h_2) = gh_2$. So $gh_1 = gh_2$, and by left cancellation $h_1 = h_2$. To show that ϕ is onto is easy. By definition every element of gH is of the form gh for some $h \in H$ and $\phi(h) = gh$.   □

Theorem 6.10 Lagrange. *Let G be a finite group and let H be a subgroup of G. Then $|G|/|H| = [G : H]$ is the number of distinct left cosets of H in G. In particular, the number of elements in H must divide the number of elements in G.*

PROOF.   The group G is partitioned into $[G : H]$ distinct left cosets. Each left coset has $|H|$ elements; therefore, $|G| = [G : H]|H|$.   □

Corollary 6.11. *Suppose that G is a finite group and $g \in G$. Then the order of g must divide the number of elements in G.*

Corollary 6.12. *Let $|G| = p$ with p a prime number. Then G is cyclic and any $g \in G$ such that $g \neq e$ is a generator.*

PROOF.   Let g be in G such that $g \neq e$. Then by Corollary 6.11, the order of g must divide the order of the group. Since $|\langle g \rangle| > 1$, it must be p. Hence, g generates G.   □

Corollary 6.12 suggests that groups of prime order p must somehow look like $\mathbb{Z}_p$.

Corollary 6.13. *Let H and K be subgroups of a finite group G such that $G \supset H \supset K$. Then*

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

PROOF.   Observe that

$$[G : K] = \frac{|G|}{|K|} = \frac{|G|}{|H|} \cdot \frac{|H|}{|K|} = [G : H][H : K].$$ □

Remark 6.14 The converse of Lagrange's Theorem is false. The group A_4 has order 12; however, it can be shown that it does not possess a subgroup of order 6. According to Lagrange's Theorem, subgroups of a group of order 12 can have orders of either 1, 2, 3, 4, or 6. However, we are not guaranteed that subgroups of every possible order exist. To prove that A_4 has no subgroup of order 6, we will assume that it does have such a subgroup H and show that a contradiction must occur. Since A_4 contains eight 3-cycles, we know that H must contain a 3-cycle. We will show that if H contains one 3-cycle, then it must contain more than 6 elements.

Proposition 6.15. *The group A_4 has no subgroup of order 6.*

PROOF.   Since $[A_4 : H] = 2$, there are only two cosets of H in A_4. Inasmuch as one of the cosets is H itself, right and left cosets must coincide; therefore, $gH = Hg$ or $gHg^{-1} = H$ for every $g \in A_4$. Since there are eight 3-cycles in A_4, at least one 3-cycle must be in H. Without loss of generality, assume that (123) is in H. Then $(123)^{-1} = (132)$ must also be in H. Since $ghg^{-1} \in H$ for all $g \in A_4$ and all $h \in H$ and

$$(124)(123)(124)^{-1} = (124)(123)(142) = (243)$$
$$(243)(123)(243)^{-1} = (243)(123)(234) = (142)$$

we can conclude that H must have at least seven elements

$$(1), (123), (132), (243), (243)^{-1} = (234), (142), (142)^{-1} = (124).$$

Therefore, A_4 has no subgroup of order 6.   □

In fact, we can say more about when two cycles have the same length.

Theorem 6.16. *Two cycles τ and μ in S_n have the same length if and only if there exists a $\sigma \in S_n$ such that $\mu = \sigma\tau\sigma^{-1}$.*

PROOF.   Suppose that

$$\tau = (a_1, a_2, \ldots, a_k)$$
$$\mu = (b_1, b_2, \ldots, b_k).$$

Define σ to be the permutation

$$\sigma(a_1) = b_1$$
$$\sigma(a_2) = b_2$$
$$\vdots$$
$$\sigma(a_k) = b_k.$$

Then $\mu = \sigma\tau\sigma^{-1}$.

Conversely, suppose that $\tau = (a_1, a_2, \ldots, a_k)$ is a k-cycle and $\sigma \in S_n$. If $\sigma(a_i) = b$ and $\sigma(a_{(i \bmod k)+1}) = b'$, then $\mu(b) = b'$. Hence,

$$\mu = (\sigma(a_1), \sigma(a_2), \ldots, \sigma(a_k)).$$

Since σ is one-to-one and onto, μ is a cycle of the same length as τ.   □

6.3   Fermat's and Euler's Theorems

The **_Euler ϕ-function_** is the map $\phi : \mathbb{N} \to \mathbb{N}$ defined by $\phi(n) = 1$ for $n = 1$, and, for $n > 1$, $\phi(n)$ is the number of positive integers m with $1 \leq m < n$ and $\gcd(m, n) = 1$.

From Proposition 3.4, we know that the order of $U(n)$, the group of units in $\mathbb{Z}_n$, is $\phi(n)$. For example, $|U(12)| = \phi(12) = 4$ since the numbers that are relatively prime to 12 are 1, 5, 7, and 11. For any prime p, $\phi(p) = p - 1$. We state these results in the following theorem.

Theorem 6.17. *Let $U(n)$ be the group of units in $\mathbb{Z}_n$. Then $|U(n)| = \phi(n)$.*

The following theorem is an important result in number theory, due to Leonhard Euler.

Theorem 6.18 Euler's Theorem. *Let a and n be integers such that $n > 0$ and $\gcd(a, n) = 1$. Then $a^{\phi(n)} \equiv 1 \pmod{n}$.*

PROOF.   By Theorem 6.17 the order of $U(n)$ is $\phi(n)$. Consequently, $a^{\phi(n)} = 1$ for all $a \in U(n)$; or $a^{\phi(n)} - 1$ is divisible by n. Therefore, $a^{\phi(n)} \equiv 1 \pmod{n}$. $\square$

If we consider the special case of Euler's Theorem in which $n = p$ is prime and recall that $\phi(p) = p - 1$, we obtain the following result, due to Pierre de Fermat.

Theorem 6.19 Fermat's Little Theorem. *Let p be any prime number and suppose that $p \nmid a$ (p does not divide a). Then*

$$a^{p-1} \equiv 1 \pmod{p}.$$

Furthermore, for any integer b, $b^p \equiv b \pmod{p}$.

Historical Note

Joseph-Louis Lagrange (1736–1813), born in Turin, Italy, was of French and Italian descent. His talent for mathematics became apparent at an early age. Leonhard Euler recognized Lagrange's abilities when Lagrange, who was only 19, communicated to Euler some work that he had done in the calculus of variations. That year he was also named a professor at the Royal Artillery School in Turin. At the age of 23 he joined the Berlin Academy. Frederick the Great had written to Lagrange proclaiming that the "greatest king in Europe" should have the "greatest mathematician in Europe" at his court. For 20 years Lagrange held the position vacated by his mentor, Euler. His works include contributions to number theory, group theory, physics and mechanics, the calculus of variations, the theory of equations, and differential equations. Along with Laplace and Lavoisier, Lagrange was one of the people responsible for designing the metric system. During his life Lagrange profoundly influenced the development of mathematics, leaving much to the next generation of mathematicians in the form of examples and new problems to be solved.

6.4   Exercises

1.   Suppose that G is a finite group with an element g of order 5 and an element h of order 7. Why must $|G| \geq 35$?

2.   Suppose that G is a finite group with 60 elements. What are the orders of possible subgroups of G?

3.   Prove or disprove: Every subgroup of the integers has finite index.

4.   Prove or disprove: Every subgroup of the integers has finite order.

5.   List the left and right cosets of the subgroups in each of the following.

(a) $\langle 8 \rangle$ in $\mathbb{Z}_{24}$ (e) A_n in S_n

(b) $\langle 3 \rangle$ in $U(8)$ (f) D_4 in S_4

(c) $3\mathbb{Z}$ in $\mathbb{Z}$ (g) $\mathbb{T}$ in $\mathbb{C}^*$

(d) A_4 in S_4 (h) $H = \{(1), (123), (132)\}$ in S_4

6.   Describe the left cosets of $SL_2(\mathbb{R})$ in $GL_2(\mathbb{R})$. What is the index of $SL_2(\mathbb{R})$ in $GL_2(\mathbb{R})$?

7. Verify Euler's Theorem for $n = 15$ and $a = 4$.

8. Use Fermat's Little Theorem to show that if $p = 4n + 3$ is prime, there is no solution to the equation $x^2 \equiv -1 \pmod{p}$.

9. Show that the integers have infinite index in the additive group of rational numbers.

10. Show that the additive group of real numbers has infinite index in the additive group of the complex numbers.

11. Let H be a subgroup of a group G and suppose that $g_1, g_2 \in G$. Prove that the following conditions are equivalent.

(a) $g_1 H = g_2 H$

(b) $H g_1^{-1} = H g_2^{-1}$

(c) $g_1 H \subset g_2 H$

(d) $g_2 \in g_1 H$

(e) $g_1^{-1} g_2 \in H$

12. If $ghg^{-1} \in H$ for all $g \in G$ and $h \in H$, show that right cosets are identical to left cosets. That is, show that $gH = Hg$ for all $g \in G$.

13. What fails in the proof of Theorem 6.8 if $\phi : \mathcal{L}_H \to \mathcal{R}_H$ is defined by $\phi(gH) = Hg$?

14. Suppose that $g^n = e$. Show that the order of g divides n.

15. Show that any two permutations $\alpha, \beta \in S_n$ have the same cycle structure if and only if there exists a permutation γ such that $\beta = \gamma \alpha \gamma^{-1}$. If $\beta = \gamma \alpha \gamma^{-1}$ for some $\gamma \in S_n$, then α and β are **_conjugate_**.

16. If $|G| = 2n$, prove that the number of elements of order 2 is odd. Use this result to show that G must contain a subgroup of order 2.

17. Suppose that $[G : H] = 2$. If a and b are not in H, show that $ab \in H$.

18. If $[G : H] = 2$, prove that $gH = Hg$.

19. Let H and K be subgroups of a group G. Prove that $gH \cap gK$ is a coset of $H \cap K$ in G.

20. Let H and K be subgroups of a group G. Define a relation $\sim$ on G by $a \sim b$ if there exists an $h \in H$ and a $k \in K$ such that $hak = b$. Show that this relation is an equivalence relation. The corresponding equivalence classes are called **_double cosets_**. Compute the double cosets of $H = \{(1), (123), (132)\}$ in A_4.

21. Let G be a cyclic group of order n. Show that there are exactly $\phi(n)$ generators for G.

22. Let $n = p_1^{e_1} p_2^{e_2} \cdots p_k^{e_k}$, where $p_1, p_2, \ldots, p_k$ are distinct primes. Prove that

$$\phi(n) = n \left(1 - \frac{1}{p_1}\right) \left(1 - \frac{1}{p_2}\right) \cdots \left(1 - \frac{1}{p_k}\right).$$

23. Show that

$$n = \sum_{d \mid n} \phi(d)$$

for all positive integers n.

6.5   Sage

Sage can create all of the cosets of a subgroup, and all of the subgroups of a group. While these methods can be somewhat slow, they are in many, many ways much better than experimenting with pencil and paper, and can greatly assist us in understanding the structure of finite groups.

Cosets

Sage will create all the right (or left) cosets of a subgroup. Written mathematically, cosets are sets, and the order of the elements within the set is irrelevant. With Sage, lists are more natural, and here it is to our advantage.

Sage creates the cosets of a subgroup as a list of lists. Each inner list is a single coset. The first coset is always the coset that is the subgroup itself, and the first element of this coset is the identity. Each of the other cosets can be construed to have their first element as their representative, and if you use this element as the representative, the elements of the coset are in the same order they would be created by multiplying this representative by the elements of the first coset (the subgroup).

The keyword `side` can be `'right'` or `'left'`, and if not given, then the default is right cosets. The options refer to which side of the product has the representative. Notice that now Sage's results will be "backwards" compared with the text. Here is Example 6.2 reprised, but in a slightly different order.

```
G = SymmetricGroup(3)
a = G("(1,2)")
H = G.subgroup([a])
rc = G.cosets(H, side='right'); rc
```

```
[[(), (1,2)], [(2,3), (1,3,2)], [(1,2,3), (1,3)]]
```

```
lc = G.cosets(H, side='left'); lc
```

```
[[(), (1,2)], [(2,3), (1,2,3)], [(1,3,2), (1,3)]]
```

So if we work our way through the brackets carefully we can see the difference between the right cosets and the left cosets. Compare these cosets with the ones in the text and see that left and right are reversed. Shouldn't be a problem — just keep it in mind.

```
G = SymmetricGroup(3)
b = G("(1,2,3)")
H = G.subgroup([b])
rc = G.cosets(H, side='right'); rc
```

```
[[(), (1,2,3), (1,3,2)], [(2,3), (1,3), (1,2)]]
```

```
lc = G.cosets(H, side='left'); lc
```

```
[[(), (1,2,3), (1,3,2)], [(2,3), (1,2), (1,3)]]
```

If we study the bracketing, we can see that the left and right cosets are equal. Let's see what Sage thinks:

```
rc == lc
```

```
False
```

Mathematically, we need sets, but Sage is working with ordered lists, and the order matters. However, if we know our lists do not have duplicates (the `.cosets()` method will never produce duplicates) then we can sort the lists and a test for equality will perform as expected. The elements of a permutation group have an ordering defined for them — it is not so important *what* this is, just that *some* ordering is defined. The `sorted()` function will take any list and return a sorted version. So for each list of cosets, we will sort the individual cosets and then sort the list of sorted cosets. This is a typical maneuver, though a bit complicated with the nested lists.

```
rc_sorted = sorted([sorted(coset) for coset in rc])
rc_sorted
```

```
[[(), (1,2,3), (1,3,2)], [(2,3), (1,2), (1,3)]]
```

```
lc_sorted = sorted([sorted(coset) for coset in lc])
lc_sorted
```

```
[[(), (1,2,3), (1,3,2)], [(2,3), (1,2), (1,3)]]
```

```
rc_sorted == lc_sorted
```

```
True
```

The list of all cosets can be quite long (it will include every element of the group) and can take a few seconds to complete, even for small groups. There are more sophisticated, and faster, ways to study cosets (such as just using their representatives), but to understand these techniques you also need to understand more theory.

Subgroups

Sage can compute all of the subgroups of a group. This can produce even more output than the coset method and can sometimes take much longer, depending on the structure of the group. The list is in order of the size of the subgroups, with smallest first. As a demonstration we will first compute and list all of the subgroups of a small group, and then extract just one of these subgroups from the list for some futher study.

```
G = SymmetricGroup(3)
sg = G.subgroups(); sg
```

```
[Subgroup of (Symmetric group of order 3! as a permutation group)
    generated by [()],
 Subgroup of (Symmetric group of order 3! as a permutation group)
    generated by [(2,3)],
 Subgroup of (Symmetric group of order 3! as a permutation group)
    generated by [(1,2)],
 Subgroup of (Symmetric group of order 3! as a permutation group)
    generated by [(1,3)],
 Subgroup of (Symmetric group of order 3! as a permutation group)
    generated by [(1,2,3)],
 Subgroup of (Symmetric group of order 3! as a permutation group)
    generated by [(2,3), (1,2,3)]]
```

```
H = sg[4]; H
```

```
Subgroup of (Symmetric group of order 3! as a permutation group)
    generated by [(1,2,3)]
```

```
H.order()
```

3

```
H.list()
```

```
[(), (1,2,3), (1,3,2)]
```

```
H.is_cyclic()
```

True

The output of the `.subgroups()` method can be voluminous, so sometimes we are interested in properties of specific subgroups (as in the previous example) or broader questions of the group's "subgroup structure." Here we expand on Corollary 6.15. Notice that just because Sage does not *compute* a subgroup of order 6 in A_4, this is no substitute whatsoever for a *proof* such as given for the corollary. But the computational result emboldens us to search for the theoretical result with confidence.

```
G = AlternatingGroup(4)
sg = G.subgroups()
[H.order() for H in sg]
```

```
[1, 2, 2, 2, 3, 3, 3, 3, 4, 12]
```

So we see no subgroup of order 6 in the list of subgroups of A_4. Notice how Lagrange's Theorem (Theorem 6.10) is in evidence — all the subgroup orders divide 12, the order of A_4. Be patient, the next subgroup computation may take a while.

```
G = SymmetricGroup(4)
sg = G.subgroups()
[H.order() for H in sg]
```

```
[1, 2, 2, 2, 2, 2, 2, 2, 2, 2, 3, 3, 3, 3, 4, 4, 4, 4, 4, 4, 4,
 6, 6, 6, 6, 8, 8, 8, 12, 24]
```

Again, note Lagrange's Theorem in action. But more interestingly, S_4 has a subgroup of order 6. Four of them, to be precise. These four subgroups of order 6 are similar to each other, can you describe them simply (*before* digging into the `sg` list for more information)? If you were curious how many subgroups S_4 has, you could simply count the number of subgroups in the `sg` list. The `len()` function does this for *any* list and is often an easy way to count things.

```
len(sg)
```

30

Subgroups of Cyclic Groups

Now that we are more familiar with permutation groups, and know about the `.subgroups()` method, we can revisit an idea from Chapter 4. The subgroups of a cyclic group are always cyclic, but how many are there and what are their orders?

```
G = CyclicPermutationGroup(20)
[H.order() for H in G.subgroups()]
```

```
[1, 2, 4, 5, 10, 20]
```

```
G = CyclicPermutationGroup(19)
[H.order() for H in G.subgroups()]
```

```
[1, 19]
```

We could do this all day, but you have Sage at your disposal, so vary the order of `G` by changing `n` and study the output across many runs. Maybe try a cyclic group of order 24 and compare with the symmetric group S_4 (above) which also has order 24. Do you feel a conjecture coming on?

```
n = 8
G = CyclicPermutationGroup(n)
[H.order() for H in G.subgroups()]
```

```
[1, 2, 4, 8]
```

Euler Phi Function

To add to our number-theoretic functions from Chapter 2, we note that Sage makes the Euler ϕ-function available as the function `euler_phi()`.

```
euler_phi(345)
```

```
176
```

Here's an interesting experiment that you can try running several times.

```
m = random_prime(10000)
n = random_prime(10000)
m, n, euler_phi(m*n) == euler_phi(m)*euler_phi(n)
```

```
(5881, 1277, True)
```

Feel another conjecture coming on? Can you generalize this result?

6.6 Sage Exercises

The following exercises are less about cosets and subgroups, and more about using Sage as an experimental tool. They are designed to help you become both more efficient, and more expressive, as you write commands in Sage. We will have many opportunities to work with cosets and subgroups in the coming chapters. These exercises do not contain much guidance, and get more challenging as they go. They are designed to explore, or confirm, results presented in this chapter or earlier chapters.

Important: You should answer each of the last three problems with a single (complicated) line of Sage that concludes by outputting `True`. A "single line" means you will have several Sage commands packaged up together in complicated ways. It does not mean several Sage commands seperated by semi-colons and typed in on a single line. Be sure include some intermediate steps used in building up your solution, but using smaller ranges of values so as to not overwhelm the reader with lots of output. This will help you, and the grader of your work, have some confidence that the final version is correct.

When you check integers below for divisibility, remember that `range()` produces plain integers, which are quite simple in their functionality. The `srange()` command produces Sage integers, which have many more capabilities. (See the last exercise for an example.) And remember that a list comprehension is a very compact way to examine many possibilities at once.**1.** Use `.subgroups()` to find an example of a group G and an integer m, so

that (a) m divides the order of G, and (b) G has no subgroup of order m. (Do not use the group A_4 for G, since this is in the text.) Provide a single line of Sage code that has all the logic to produce the desired m as its output. (You can give your group a simple name on a prior line and then just reference the group by name.) Here is a very simple example that might help you structure your answer.

```
a = 5
b = 10
c = 6
d = 13
a.divides(b)
```

True

```
not (b in [c,d])
```

True

```
a.divides(b) and not (b in [c,d])
```

True

2.   Verify the truth of Fermat's Little Theorem (either variant) using the composite number $391 = 17 \cdot 23$ as the choice of the base (either a or b), and for p assuming the value of every prime number between 100 and 1000.

Build up a solution slowly — make a list of powers (start with just a few primes), then make a list of powers reduced by modular arithmetic, then a list of comparisons with the predicted value, then a check on all these logical values resulting from the comparisons. This is a useful strategy for many similar problems. Eventually you will write a single line that performs the verification by eventually printing out True. Here are some more hints about useful functions.

```
a = 20
b = 6
a.mod(b)
```

2

```
prime_range(50, 100)
```

[53, 59, 61, 67, 71, 73, 79, 83, 89, 97]

```
all([True, True, True, True])
```

True

```
all([True, True, False, True])
```

False

3.   Verify that the group of units mod n has order $n - 1$ when n is prime, again for all primes between 100 and 1000. As before, your output should be simply True, just once, indicating that the statement about the order is true for all the primes examined. As before, build up your solution slowly, and with a smaller range of primes in the beginning. Express your answer as a single line of Sage code.

4. Verify Euler's Theorem for all values of $0 < n < 100$ and for $1 \leq a \leq n$. This will require nested `for` statements with a conditional. Again, here is a small example that might be helpful for constructing your one line of Sage code. Note the use of `srange()` in this example.

```
[a/b for a in srange(9) for b in srange(1,a) if gcd(a,b)==1]
```

```
[2, 3, 3/2, 4, 4/3, 5, 5/2, 5/3, 5/4, 6, 6/5,
 7, 7/2, 7/3, 7/4, 7/5, 7/6, 8, 8/3, 8/5, 8/7]
```

5. The symmetric group on 7 symbols, S_7, has $7! = 5040$ elements. Consider the following questions without employing Sage, based on what we know about orders of elements of permutation groups (Exercise 5.3.13).

- What is the maximum possible order?

- How many elements are there of order 10?

- How many elements are there of order 1?

- How many elements are there of order 2?

- What is the smallest positive integer for which there is no element with that order?

These questions will be easier if you are familiar with using binomial coefficients for counting in similarly complex situations. But either way, give some serious thought to each question (and maybe a few of your own) before firing up Sage.

Now, compute how many elements there are of each order using the `.order()` method, and then embed this into a list comprehension which creates a single list of these counts. You can check your work (or check Sage) by wrapping this list in `sum()` and hopefully getting 5040.

Comment on the process of studying these questions first without any computational aid, and then again with Sage. For which values of n do you think Sage would be too slow and your mind quicker?

7

Introduction to Cryptography

Cryptography is the study of sending and receiving secret messages. The aim of cryptography is to send messages across a channel so that only the intended recipient of the message can read it. In addition, when a message is received, the recipient usually requires some assurance that the message is authentic; that is, that it has not been sent by someone who is trying to deceive the recipient. Modern cryptography is heavily dependent on abstract algebra and number theory.

The message to be sent is called the ***plaintext*** message. The disguised message is called the ***ciphertext***. The plaintext and the ciphertext are both written in an ***alphabet***, consisting of ***letters*** or ***characters***. Characters can include not only the familiar alphabetic characters A, ..., Z and a, ..., z but also digits, punctuation marks, and blanks. A ***cryptosystem***, or ***cipher***, has two parts: ***encryption***, the process of transforming a plaintext message to a ciphertext message, and ***decryption***, the reverse transformation of changing a ciphertext message into a plaintext message.

There are many different families of cryptosystems, each distinguished by a particular encryption algorithm. Cryptosystems in a specified cryptographic family are distinguished from one another by a parameter to the encryption function called a ***key***. A classical cryptosystem has a single key, which must be kept secret, known only to the sender and the receiver of the message. If person A wishes to send secret messages to two different people B and C, and does not wish to have B understand C's messages or vice versa, A must use two separate keys, so one cryptosystem is used for exchanging messages with B, and another is used for exchanging messages with C.

Systems that use two separate keys, one for encoding and another for decoding, are called ***public key cryptosystems***. Since knowledge of the encoding key does not allow anyone to guess at the decoding key, the encoding key can be made public. A public key cryptosystem allows A and B to send messages to C using the same encoding key. Anyone is capable of encoding a message to be sent to C, but only C knows how to decode such a message.

7.1 Private Key Cryptography

In ***single*** or ***private key cryptosystems*** the same key is used for both encrypting and decrypting messages. To encrypt a plaintext message, we apply to the message some function which is kept secret, say f. This function will yield an encrypted message. Given the encrypted form of the message, we can recover the original message by applying the inverse transformation f^{-1}. The transformation f must be relatively easy to compute, as must f^{-1}; however, f must be extremely difficult to guess from available examples of coded messages.

Example 7.1. One of the first and most famous private key cryptosystems was the shift code used by Julius Caesar. We first digitize the alphabet by letting $A = 00, B = 01, \ldots, Z = 25$. The encoding function will be

$$f(p) = p + 3 \bmod 26;$$

that is, $A \mapsto D, B \mapsto E, \ldots, Z \mapsto C$. The decoding function is then

$$f^{-1}(p) = p - 3 \bmod 26 = p + 23 \bmod 26.$$

Suppose we receive the encoded message DOJHEUD. To decode this message, we first digitize it:

$$3, 14, 9, 7, 4, 20, 3.$$

Next we apply the inverse transformation to get

$$0, 11, 6, 4, 1, 17, 0,$$

or ALGEBRA. Notice here that there is nothing special about either of the numbers 3 or 26. We could have used a larger alphabet or a different shift.

Cryptanalysis is concerned with deciphering a received or intercepted message. Methods from probability and statistics are great aids in deciphering an intercepted message; for example, the frequency analysis of the characters appearing in the intercepted message often makes its decryption possible.

Example 7.2. Suppose we receive a message that we know was encrypted by using a shift transformation on single letters of the 26-letter alphabet. To find out exactly what the shift transformation was, we must compute b in the equation $f(p) = p + b \bmod 26$. We can do this using frequency analysis. The letter $E = 04$ is the most commonly occurring letter in the English language. Suppose that $S = 18$ is the most commonly occurring letter in the ciphertext. Then we have good reason to suspect that $18 = 4 + b \bmod 26$, or $b = 14$. Therefore, the most likely encrypting function is

$$f(p) = p + 14 \bmod 26.$$

The corresponding decrypting function is

$$f^{-1}(p) = p + 12 \bmod 26.$$

It is now easy to determine whether or not our guess is correct.

Simple shift codes are examples of **monoalphabetic cryptosystems**. In these ciphers a character in the enciphered message represents exactly one character in the original message. Such cryptosystems are not very sophisticated and are quite easy to break. In fact, in a simple shift as described in Example 7.1, there are only 26 possible keys. It would be quite easy to try them all rather than to use frequency analysis.

Let us investigate a slightly more sophisticated cryptosystem. Suppose that the encoding function is given by

$$f(p) = ap + b \bmod 26.$$

We first need to find out when a decoding function f^{-1} exists. Such a decoding function exists when we can solve the equation

$$c = ap + b \bmod 26$$

for p. By Proposition 3.4, this is possible exactly when a has an inverse or, equivalently, when $\gcd(a, 26) = 1$. In this case

$$f^{-1}(p) = a^{-1}p - a^{-1}b \bmod 26.$$

Such a cryptosystem is called an **affine cryptosystem**.

Example 7.3. Let us consider the affine cryptosystem $f(p) = ap + b \bmod 26$. For this cryptosystem to work we must choose an $a \in \mathbb{Z}_{26}$ that is invertible. This is only possible if $\gcd(a, 26) = 1$. Recognizing this fact, we will let $a = 5$ since $\gcd(5, 26) = 1$. It is easy to see that $a^{-1} = 21$. Therefore, we can take our encryption function to be $f(p) = 5p + 3 \bmod 26$. Thus, ALGEBRA is encoded as $3, 6, 7, 23, 8, 10, 3$, or DGHXIKD. The decryption function will be
$$f^{-1}(p) = 21p - 21 \cdot 3 \bmod 26 = 21p + 15 \bmod 26.$$

A cryptosystem would be more secure if a ciphertext letter could represent more than one plaintext letter. To give an example of this type of cryptosystem, called a **polyalphabetic cryptosystem**, we will generalize affine codes by using matrices. The idea works roughly the same as before; however, instead of encrypting one letter at a time we will encrypt pairs of letters. We can store a pair of letters p_1 and p_2 in a vector

$$\mathbf{p} = \begin{pmatrix} p_1 \\ p_2 \end{pmatrix}.$$

Let A be a 2×2 invertible matrix with entries in $\mathbb{Z}_{26}$. We can define an encoding function by
$$f(\mathbf{p}) = A\mathbf{p} + \mathbf{b},$$

where $\mathbf{b}$ is a fixed column vector and matrix operations are performed in $\mathbb{Z}_{26}$. The decoding function must be
$$f^{-1}(\mathbf{p}) = A^{-1}\mathbf{p} - A^{-1}\mathbf{b}.$$

Example 7.4. Suppose that we wish to encode the word HELP. The corresponding digit string is $7, 4, 11, 15$. If
$$A = \begin{pmatrix} 3 & 5 \\ 1 & 2 \end{pmatrix},$$

then
$$A^{-1} = \begin{pmatrix} 2 & 21 \\ 25 & 3 \end{pmatrix}.$$

If $\mathbf{b} = (2, 2)^{\mathrm{t}}$, then our message is encrypted as RRGR. The encrypted letter R represents more than one plaintext letter.

Frequency analysis can still be performed on a polyalphabetic cryptosystem, because we have a good understanding of how pairs of letters appear in the English language. The pair *th* appears quite often; the pair *qz* never appears. To avoid decryption by a third party, we must use a larger matrix than the one we used in Example 7.4.

7.2 Public Key Cryptography

If traditional cryptosystems are used, anyone who knows enough to encode a message will also know enough to decode an intercepted message. In 1976, W. Diffie and M. Hellman proposed public key cryptography, which is based on the observation that the encryption and

decryption procedures need not have the same key. This removes the requirement that the encoding key be kept secret. The encoding function f must be relatively easy to compute, but f^{-1} must be extremely difficult to compute without some additional information, so that someone who knows only the encrypting key cannot find the decrypting key without prohibitive computation. It is interesting to note that to date, no system has been proposed that has been proven to be "one-way;" that is, for any existing public key cryptosystem, it has never been shown to be computationally prohibitive to decode messages with only knowledge of the encoding key.

The RSA Cryptosystem

The RSA cryptosystem introduced by R. Rivest, A. Shamir, and L. Adleman in 1978, is based on the difficulty of factoring large numbers. Though it is not a difficult task to find two large random primes and multiply them together, factoring a 150-digit number that is the product of two large primes would take 100 million computers operating at 10 million instructions per second about 50 million years under the fastest algorithms available in the early 1990s. Although the algorithms have improved, factoring a number that is a product of two large primes is still computationally prohibitive.

The RSA cryptosystem works as follows. Suppose that we choose two random 150-digit prime numbers p and q. Next, we compute the product $n = pq$ and also compute $\phi(n) = m = (p-1)(q-1)$, where ϕ is the Euler ϕ-function. Now we start choosing random integers E until we find one that is relatively prime to m; that is, we choose E such that $\gcd(E, m) = 1$. Using the Euclidean algorithm, we can find a number D such that $DE \equiv 1 \pmod{m}$. The numbers n and E are now made public.

Suppose now that person B (Bob) wishes to send person A (Alice) a message over a public line. Since E and n are known to everyone, anyone can encode messages. Bob first digitizes the message according to some scheme, say $A = 00, B = 02, \ldots, Z = 25$. If necessary, he will break the message into pieces such that each piece is a positive integer less than n. Suppose x is one of the pieces. Bob forms the number $y = x^E \mod n$ and sends y to Alice. For Alice to recover x, she need only compute $x = y^D \mod n$. Only Alice knows D.

Example 7.5. Before exploring the theory behind the RSA cryptosystem or attempting to use large integers, we will use some small integers just to see that the system does indeed work. Suppose that we wish to send some message, which when digitized is 25. Let $p = 23$ and $q = 29$. Then

$$n = pq = 667$$

and

$$\phi(n) = m = (p-1)(q-1) = 616.$$

We can let $E = 487$, since $\gcd(616, 487) = 1$. The encoded message is computed to be

$$25^{487} \mod 667 = 169.$$

This computation can be reasonably done by using the method of repeated squares as described in Chapter 4. Using the Euclidean algorithm, we determine that $191E = 1 + 151m$; therefore, the decrypting key is $(n, D) = (667, 191)$. We can recover the original message by calculating

$$169^{191} \mod 667 = 25.$$

Now let us examine why the RSA cryptosystem works. We know that $DE \equiv 1 \pmod{m}$; hence, there exists a k such that

$$DE = km + 1 = k\phi(n) + 1.$$

There are two cases to consider. In the first case assume that $\gcd(x, n) = 1$. Then by Theorem 6.18,

$$y^D = (x^E)^D = x^{DE} = x^{km+1} = (x^{\phi(n)})^k x = (1)^k x = x \bmod n.$$

So we see that Alice recovers the original message x when she computes $y^D \bmod n$.

For the other case, assume that $\gcd(x, n) \neq 1$. Since $n = pq$ and $x < n$, we know x is a multiple of p or a multiple of q, but not both. We will describe the first possibility only, since the second is entirely similar. There is then an integer r, with $r < q$ and $x = rp$. Note that we have $\gcd(x, q) = 1$ and that $m = \phi(n) = (p-1)(q-1) = \phi(p)\phi(q)$. Then, using Theorem 6.18, but now mod q,

$$x^{km} = x^{k\phi(p)\phi(q)} = (x^{\phi(q)})^{k\phi(p)} = (1)^{k\phi(p)} = 1 \bmod q.$$

So there is an integer t such that $x^{km} = 1 + tq$. Thus, Alice also recovers the message in this case,

$$y^D = x^{km+1} = x^{km}x = (1 + tq)x = x + tq(rp) = x + trn = x \bmod n.$$

We can now ask how one would go about breaking the RSA cryptosystem. To find D given n and E, we simply need to factor n and solve for D by using the Euclidean algorithm. If we had known that $667 = 23 \cdot 29$ in Example 7.5, we could have recovered D.

Message Verification

There is a problem of message verification in public key cryptosystems. Since the encoding key is public knowledge, anyone has the ability to send an encoded message. If Alice receives a message from Bob, she would like to be able to verify that it was Bob who actually sent the message. Suppose that Bob's encrypting key is (n', E') and his decrypting key is (n', D'). Also, suppose that Alice's encrypting key is (n, E) and her decrypting key is (n, D). Since encryption keys are public information, they can exchange coded messages at their convenience. Bob wishes to assure Alice that the message he is sending is authentic. Before Bob sends the message x to Alice, he decrypts x with his own key:

$$x' = x^{D'} \bmod n'.$$

Anyone can change x' back to x just by encryption, but only Bob has the ability to form x'. Now Bob encrypts x' with Alice's encryption key to form

$$y' = x'^E \bmod n,$$

a message that only Alice can decode. Alice decodes the message and then encodes the result with Bob's key to read the original message, a message that could have only been sent by Bob.

Historical Note

Encrypting secret messages goes as far back as ancient Greece and Rome. As we know, Julius Caesar used a simple shift code to send and receive messages. However, the formal study of encoding and decoding messages probably began with the Arabs in the 1400s. In the fifteenth and sixteenth centuries mathematicians such as Alberti and Viete discovered that monoalphabetic cryptosystems offered no real security. In the 1800s, F. W. Kasiski established methods for breaking ciphers in which a ciphertext letter can represent more

than one plaintext letter, if the same key was used several times. This discovery led to the use of cryptosystems with keys that were used only a single time. Cryptography was placed on firm mathematical foundations by such people as W. Friedman and L. Hill in the early part of the twentieth century.

The period after World War I saw the development of special-purpose machines for encrypting and decrypting messages, and mathematicians were very active in cryptography during World War II. Efforts to penetrate the cryptosystems of the Axis nations were organized in England and in the United States by such notable mathematicians as Alan Turing and A. A. Albert. The Allies gained a tremendous advantage in World War II by breaking the ciphers produced by the German Enigma machine and the Japanese Purple ciphers.

By the 1970s, interest in commercial cryptography had begun to take hold. There was a growing need to protect banking transactions, computer data, and electronic mail. In the early 1970s, IBM developed and implemented LUZIFER, the forerunner of the National Bureau of Standards' Data Encryption Standard (DES).

The concept of a public key cryptosystem, due to Diffie and Hellman, is very recent (1976). It was further developed by Rivest, Shamir, and Adleman with the RSA cryptosystem (1978). It is not known how secure any of these systems are. The trapdoor knapsack cryptosystem, developed by Merkle and Hellman, has been broken. It is still an open question whether or not the RSA system can be broken. In 1991, RSA Laboratories published a list of semiprimes (numbers with exactly two prime factors) with a cash prize for whoever was able to provide a factorization (http://www.emc.com/emc-plus/rsa-labs/historical/the-rsa-challenge-numbers.htm). Although the challenge ended in 2007, many of these numbers have not yet been factored.

There been a great deal of controversy about research in cryptography and cryptography itself. In 1929, when Henry Stimson, Secretary of State under Herbert Hoover, dismissed the Black Chamber (the State Department's cryptography division) on the ethical grounds that "gentlemen do not read each other's mail." During the last two decades of the twentieth century, the National Security Agency wanted to keep information about cryptography secret, whereas the academic community fought for the right to publish basic research. Currently, research in mathematical cryptography and computational number theory is very active, and mathematicians are free to publish their results in these areas.

7.3 Exercises

1. Encode `IXLOVEXMATH` using the cryptosystem in Example 7.1.

2. Decode `ZLOOA WKLVA EHARQ WKHA ILQDO`, which was encoded using the cryptosystem in Example 7.1.

3. Assuming that monoalphabetic code was used to encode the following secret message, what was the original message?

```
APHUO EGEHP PEXOV FKEUH CKVUE CHKVE APHUO
EGEHU EXOVL EXDKT VGEFT EHFKE UHCKF TZEXO
VEZDT TVKUE XOVKV ENOHK ZFTEH TEHKQ LEROF
PVEHP PEXOV ERYKP GERYT GVKEG XDRTE RGAGA
```

What is the significance of this message in the history of cryptography?

4. What is the total number of possible monoalphabetic cryptosystems? How secure are such cryptosystems?

5. Prove that a 2×2 matrix A with entries in $\mathbb{Z}_{26}$ is invertible if and only if $\gcd(\det(A), 26) = 1$.

6. Given the matrix

$$A = \begin{pmatrix} 3 & 4 \\ 2 & 3 \end{pmatrix},$$

use the encryption function $f(\mathbf{p}) = A\mathbf{p} + \mathbf{b}$ to encode the message CRYPTOLOGY, where $\mathbf{b} = (2,5)^{\mathrm{t}}$. What is the decoding function?

7. Encrypt each of the following RSA messages x so that x is divided into blocks of integers of length 2; that is, if $x = 142528$, encode 14, 25, and 28 separately.

(a) $n = 3551, E = 629, x = 31$

(b) $n = 2257, E = 47, x = 23$

(c) $n = 120979, E = 13251, x = 142371$

(d) $n = 45629, E = 781, x = 231561$

8. Compute the decoding key D for each of the encoding keys in Exercise 7.3.7.

9. Decrypt each of the following RSA messages y.

(a) $n = 3551, D = 1997, y = 2791$

(b) $n = 5893, D = 81, y = 34$

(c) $n = 120979, D = 27331, y = 112135$

(d) $n = 79403, D = 671, y = 129381$

10. For each of the following encryption keys (n, E) in the RSA cryptosystem, compute D.

(a) $(n, E) = (451, 231)$

(b) $(n, E) = (3053, 1921)$

(c) $(n, E) = (37986733, 12371)$

(d) $(n, E) = (16394854313, 34578451)$

11. Encrypted messages are often divided into blocks of n letters. A message such as THE WORLD WONDERS WHY might be encrypted as JIW OCFRJ LPOEVYQ IOC but sent as JIW OCF RJL POE VYQ IOC. What are the advantages of using blocks of n letters?

12. Find integers n, E, and X such that

$$X^E \equiv X \pmod{n}.$$

Is this a potential problem in the RSA cryptosystem?

13. Every person in the class should construct an RSA cryptosystem using primes that are 10 to 15 digits long. Hand in (n, E) and an encoded message. Keep D secret. See if you can break one another's codes.

7.4 Additional Exercises: Primality and Factoring

In the RSA cryptosystem it is important to be able to find large prime numbers easily. Also, this cryptosystem is not secure if we can factor a composite number that is the product of two large primes. The solutions to both of these problems are quite easy. To find out if a number n is prime or to factor n, we can use trial division. We simply divide n by $d = 2, 3, \ldots, \sqrt{n}$. Either a factorization will be obtained, or n is prime if no d divides n. The problem is that such a computation is prohibitively time-consuming if n is very large.**1.** A better algorithm for factoring odd positive integers is ***Fermat's factorization algorithm***.

(a) Let $n = ab$ be an odd composite number. Prove that n can be written as the difference of two perfect squares:
$$n = x^2 - y^2 = (x - y)(x + y).$$

Consequently, a positive odd integer can be factored exactly when we can find integers x and y such that $n = x^2 - y^2$.

(b) Write a program to implement the following factorization algorithm based on the observation in part (a). The expression `ceiling(sqrt(n))` means the smallest integer greater than or equal to the square root of n. Write another program to do factorization using trial division and compare the speed of the two algorithms. Which algorithm is faster and why?

```
x := ceiling(sqrt(n))
y := 1

1 : while x^2 - y^2 > n do
      y := y + 1

if x^2 - y^2 < n then
      x := x + 1
      y := 1
      goto 1
else if x^2 - y^2 = 0 then
      a := x - y
      b := x + y
      write n = a * b
```

2. **Primality Testing.** Recall Fermat's Little Theorem from Chapter 6. Let p be prime with $\gcd(a, p) = 1$. Then $a^{p-1} \equiv 1 \pmod{p}$. We can use Fermat's Little Theorem as a screening test for primes. For example, 15 cannot be prime since
$$2^{15-1} \equiv 2^{14} \equiv 4 \pmod{15}.$$

However, 17 is a potential prime since
$$2^{17-1} \equiv 2^{16} \equiv 1 \pmod{17}.$$

We say that an odd composite number n is a ***pseudoprime*** if
$$2^{n-1} \equiv 1 \pmod{n}.$$

Which of the following numbers are primes and which are pseudoprimes?

(a) 342   (c) 601   (e) 771

(b) 811   (d) 561   (f) 631

3.   Let n be an odd composite number and b be a positive integer such that $\gcd(b, n) = 1$. If $b^{n-1} \equiv 1 \pmod{n}$, then n is a **pseudoprime base** b. Show that 341 is a pseudoprime base 2 but not a pseudoprime base 3.

4.   Write a program to determine all primes less than 2000 using trial division. Write a second program that will determine all numbers less than 2000 that are either primes or pseudoprimes. Compare the speed of the two programs. How many pseudoprimes are there below 2000?

There exist composite numbers that are pseudoprimes for all bases to which they are relatively prime. These numbers are called **Carmichael numbers**. The first Carmichael number is $561 = 3 \cdot 11 \cdot 17$. In 1992, Alford, Granville, and Pomerance proved that there are an infinite number of Carmichael numbers [4]. However, Carmichael numbers are very rare. There are only 2163 Carmichael numbers less than 25×10^9. For more sophisticated primality tests, see [1], [6], or [7].

References and Suggested Readings

[**1**]   Bressoud, D. M. *Factorization and Primality Testing.* Springer-Verlag, New York, 1989.

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[**4**]   Granville, A. "Primality Testing and Carmichael Numbers," *Notices of the American Mathematical Society* **39**(1992), 696–700.

[**5**]   Hellman, M. E. "The Mathematics of Public Key Cryptography," *Scientific American* **241**(1979), 130–39.

[**6**]   Koblitz, N. *A Course in Number Theory and Cryptography.* 2nd ed. Springer, New York, 1994.

[**7**]   Pomerance, C., ed. "Cryptology and Computational Number Theory", *Proceedings of Symposia in Applied Mathematics* **42**(1990) American Mathematical Society, Providence, RI.

[**8**]   Rivest, R. L., Shamir, A., and Adleman, L., "A Method for Obtaining Signatures and Public-key Cryptosystems," *Comm. ACM* **21**(1978), 120–26.

7.5   Sage

Since Sage began as software to support research in number theory, we can quickly and easily demonstrate the internal workings of the RSA algorithm. Recognize that, in practice, many other details such as encoding between letters and integers, or protecting one's private key, are equally important for the security of communications. So RSA itself is just the theoretical foundation.

Constructing Keys

We will suppose that Alice wants to send a secret message to Bob, along with message verification (also known as a message with a digital signature). So we begin with the construction of key pairs (private and public) for both Alice and Bob. We first need two large primes for both individuals, and their product. In practice, values of n would have hundreds of digits, rather than just 21 as we have done here.

```
p_a = next_prime(10^10)
q_a = next_prime(p_a)
p_b = next_prime((3/2)*10^10)
q_b = next_prime(p_b)
n_a = p_a * q_a
n_b = p_b * q_b
n_a, n_b
```

(100000000520000000627, 225000000300000000091)

Computationally, the value of the Euler ϕ-function for a product of primes pq can be obtained from $(p-1)(q-1)$, but we could use Sage's built-in function just as well.

```
m_a = euler_phi(n_a)
m_b = euler_phi(n_b)
m_a, m_b
```

(100000000500000000576, 225000000270000000072)

Now we can create the encryption and decryption exponents. We choose the encryption exponent as a (small) number relatively prime to the value of m. With Sage we can factor m quickly to help us choose this value. In practice we would not want to do this computation for large values of m, so we might more easily choose "random" values and check for the first value which is relatively prime to m. The decryption exponent is the multiplicative inverse, mod m, of the encryption exponent. If you construct an improper encryption exponent (not relatively prime to m), the computation of the multiplicative inverse will fail (and Sage will tell you so). We do this twice —- for both Alice and Bob.

```
factor(m_a)
```

2^6 * 3 * 11 * 17 * 131 * 521 * 73259 * 557041

```
E_a = 5*23
D_a = inverse_mod(E_a, m_a)
D_a
```

20869565321739130555

```
factor(m_b)
```

2^3 * 3^4 * 107 * 1298027 * 2500000001

```
E_b = 7*29
D_b = inverse_mod(E_b, m_b)
D_b
```

24384236482463054195

At this stage, each individual would publish their values of n and E, while keeping D very private and secure. In practice D should be protected on the user's hard disk by a password

only the owner knows. For even greater security a person might only have two copies of their private key, one on a USB memory stick they always carry with them, and a backup in their sage deposit box. Every time the person uses D they would need to provide the password. The value of m can be discarded. For the record, here are all the keys:

```
print("Alice's_public_key,_n:", n_a, "E:", E_a)
```

Alice's_public_key,_n:_1000000000520000000627_E:_115

```
print("Alice's_private_key,_D:", D_a)
```

Alice's_private_key,_D:_20869565321739130555

```
print("Bob's_public_key,_n:", n_b, "E:", E_b)
```

Bob's_public_key,_n:_2250000000300000000091_E:_203

```
print("Bob's_private_key,_D:", D_b)
```

Bob's_private_key,_D:_24384236482463054195

Signing and Encoding a Message

Alice is going to construct a message as an English word with four letters. From these four letters we will construct a single number to represent the message in a form we can use in the RSA algorithm. The function `ord()` will convert a single letter to its ASCII code value, a number between 0 and 127. If we use these numbers as "digits" mod 128, we can be sure that Alice's four-letter word will encode to an integer less than $128^4 = 268,435,456$. The particular maximum value is not important, so long as it is smaller than our value of n since all of our subsequent arithmetic is mod n. We choose a popular four-letter word, convert to ASCII "digits" with a list comprehension, and then construct the integer from the digits with the right base. Notice how we can treat the word as a list and that the first digit in the list is in the "ones" place (we say the list is in "little-endian" order).

```
word = 'Sage'
digits = [ord(letter) for letter in word]
digits
```

[83, 97, 103, 101]

```
message = ZZ(digits, 128)
message
```

213512403

First, Alice will sign her message to provide message verification. She uses her private key for this, since this is an act that only she should be able to perform.

```
signed = power_mod(message, D_a, n_a)
signed
```

47838774644892618423

Then Alice encrypts her message so that only Bob can read it. To do this, she uses Bob's public key. Notice how she does not have to even know Bob — for example, she could have obtained Bob's public key off his web site or maybe Bob announced his public key in an advertisement in the *New York Times*.

```
encrypted = power_mod(signed, E_b, n_b)
encrypted
```

111866209291209840488

Alice's communication is now ready to travel on any communications network, no matter how insecure the network may be, and no matter how many snoops may be monitoring the network.

Decoding and Verifying a Message

Now assume that the value of `encrypted` has reached Bob. Realize that Bob may not know Alice, and realize that Bob does not even necessarily believe what he has received has genuinely originated from Alice. An adversary could be trying to confuse Bob by sending messages that claim to be from Alice. First, Bob must unwrap the encyption Alice has provided. This is an act only Bob, as the intended recipient, should be able to do. And he does it by using his private key, which only he knows, and which he has kept secure.

```
decrypted = power_mod(encrypted, D_b, n_b)
decrypted
```

47838774644892618423

Right now, this means very little to Bob. Anybody could have sent him an encoded message. However, this was a message Alice signed. Lets unwrap the message signing. Notice that this uses Alice's public key. Bob does not need to know Alice — for example, he could obtain Alice's key off her web site or maybe Alice announced her public key in an advertisement in the *New York Times*.

```
received = power_mod(decrypted, E_a, n_a)
received
```

213512403

Bob needs to transform this integer representation back to a word with letters. The `chr()` function converts ASCII code values to letters, and we use a list comprehension to do this repeatedly.

```
digits = received.digits(base=128)
letters = [chr(ascii) for ascii in digits]
letters
```

['S', 'a', 'g', 'e']

If we would like a slightly more recognizable result, we can combine the letters into a string.

```
''.join(letters)
```

'Sage'

Bob is pleased to obtain such an informative message from Alice. What would have happened if an imposter had sent a message ostensibly from Alice, or what if an adversary had intercepted Alice's original message and replaced it with a tampered message? (The latter is known as a "man in the middle" attack.)

In either case, the rogue party would not be able to duplicate Alice's first action — signing her message. If an adversary somehow signs the message, or tampers with it, the step where Bob unwraps the signing will lead to total garbage. (Try it!) Because Bob

received a legitimate word, properly capitalized, he has confidence that the message he
unsigned is the same as the message Alice signed. In practice, if Alice sent several hundred
words as her message, the odds that it will unsign as cohrent text are astronomically small.

What have we demonstrated?

1. Alice can send messages that only Bob can read.

2. Bob can receive secret messages from anybody.

3. Alice can sign messages, so that then Bob (or anybody else)knows they are genuinely
   from Alice.

Of course, without making new keys, you can reverse the roles of Alice and Bob. And if
Carol makes a key pair, she can communicate with both Alice and Bob in the same fashion.

If you want to use RSA public-key encryption seriously, investigate the open source
software GNU Privacy Guard, aka GPG, which is freely available at www.gnupg.org/. Notice
that it only makes sense to use encryption programs that allow you to look at the source
code.

7.6   Sage Exercises

1.   Construct a keypair for Alice using the first two primes greater than 10^{12}. For your
choice of E, use a single prime number and use the smallest possible choice.

Output the values of n, E, and D for Alice. Then use Sage commands to verify that
Alice's encryption and decryption keys are multiplicative inverses.

2.   Construct a keypair for Bob using the first two primes greater than $2 \cdot 10^{12}$. For your
choice of E, use a single prime number and use the smallest possible choice. Output the
values of n, E, and D for Alice.

Encode the word Math using ASCII values in the same manner as described in this section
(keep the capitalization as shown). Create a signed message of this word for communication
from Alice to Bob. Output the three integers: the message, the signed message and the
signed, encrypted message.

3.   Demonstrate how Bob converts the message received from Alice back into the word
Math. Output the value of the intermediate computations and the final human-readable
message.

4.   Create a new signed message from Alice to Bob. Simulate the message being tampered
with by adding 1 to the integer Bob receives, before he decrypts it. What result does Bob
get for the letters of the message when he decrypts and unsigns the tampered message?

5. **Classroom Exercise.**   Organize a class into several small groups. Have each group
construct key pairs with some minimum size (digits in n). Each group should keep their
private key to themselves, but make their public key available to everybody in the room.
It could be written on the board (error-prone) or maybe pasted in a public site like paste-
bin.com. Then each group can send a signed message to another group, where the groups
could be arranged logically in a circular fashion for this purpose. Of course, messages should
be posted publicly as well. Expect a success rate somewhere between 50% and 100%.

If you do not do this in class, grab a study buddy and send each other messages in the
same manner. Expect a success rate of 0%, 50% or 100%.

8

Algebraic Coding Theory

Coding theory is an application of algebra that has become increasingly important over the last several decades. When we transmit data, we are concerned about sending a message over a channel that could be affected by "noise." We wish to be able to encode and decode the information in a manner that will allow the detection, and possibly the correction, of errors caused by noise. This situation arises in many areas of communications, including radio, telephone, television, computer communications, and digital media technology. Probability, combinatorics, group theory, linear algebra, and polynomial rings over finite fields all play important roles in coding theory.

8.1   Error-Detecting and Correcting Codes

Let us examine a simple model of a communications system for transmitting and receiving coded messages (Figure 8.1).

m-digit message

Encoder

n-digit code word

Transmitter

Noise

Receiver

n-digit received word

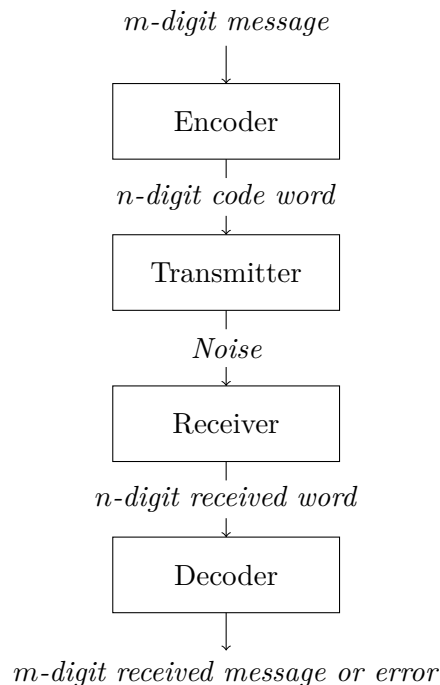| Decoder |
| :---: |

m-digit received message or error

Figure 8.1: Encoding and decoding messages

Uncoded messages may be composed of letters or characters, but typically they consist of binary m-tuples. These messages are encoded into codewords, consisting of binary n-tuples, by a device called an ***encoder***. The message is transmitted and then decoded. We

will consider the occurrence of errors during transmission. An ***error*** occurs if there is a change in one or more bits in the codeword. A ***decoding scheme*** is a method that either converts an arbitrarily received n-tuple into a meaningful decoded message or gives an error message for that n-tuple. If the received message is a codeword (one of the special n-tuples allowed to be transmitted), then the decoded message must be the unique message that was encoded into the codeword. For received non-codewords, the decoding scheme will give an error indication, or, if we are more clever, will actually try to correct the error and reconstruct the original message. Our goal is to transmit error-free messages as cheaply and quickly as possible.

Example 8.2. One possible coding scheme would be to send a message several times and to compare the received copies with one another. Suppose that the message to be encoded is a binary n-tuple $(x_1, x_2, \ldots, x_n)$. The message is encoded into a binary $3n$-tuple by simply repeating the message three times:

$$(x_1, x_2, \ldots, x_n) \mapsto (x_1, x_2, \ldots, x_n, x_1, x_2, \ldots, x_n, x_1, x_2, \ldots, x_n).$$

To decode the message, we choose as the ith digit the one that appears in the ith place in at least two of the three transmissions. For example, if the original message is (0110), then the transmitted message will be $(0110\ 0110\ 0110)$. If there is a transmission error in the fifth digit, then the received codeword will be $(0110\ 1110\ 0110)$, which will be correctly decoded as (0110).[1] This triple-repetition method will automatically detect and correct all single errors, but it is slow and inefficient: to send a message consisting of n bits, $2n$ extra bits are required, and we can only detect and correct single errors. We will see that it is possible to find an encoding scheme that will encode a message of n bits into m bits with m much smaller than $3n$.

Example 8.3. *Even parity*, a commonly used coding scheme, is much more efficient than the simple repetition scheme. The ASCII (American Standard Code for Information Interchange) coding system uses binary 8-tuples, yielding $2^8 = 256$ possible 8-tuples. However, only seven bits are needed since there are only $2^7 = 128$ ASCII characters. What can or should be done with the extra bit? Using the full eight bits, we can detect single transmission errors. For example, the ASCII codes for A, B, and C are

$$A = 65_{10} = 01000001_2,$$
$$B = 66_{10} = 01000010_2,$$
$$C = 67_{10} = 01000011_2.$$

Notice that the leftmost bit is always set to 0; that is, the 128 ASCII characters have codes

$$00000000_2 = 0_{10},$$
$$\vdots$$
$$01111111_2 = 127_{10}.$$

The bit can be used for error checking on the other seven bits. It is set to either 0 or 1 so that the total number of 1 bits in the representation of a character is even. Using even parity, the codes for A, B, and C now become

$$A = 01000001_2,$$
$$B = 01000010_2,$$

---

[1]We will adopt the convention that bits are numbered left to right in binary n-tuples.

$$C = 11000011_2.$$

Suppose an A is sent and a transmission error in the sixth bit is caused by noise over the communication channel so that (0100 0101) is received. We know an error has occurred since the received word has an odd number of 1s, and we can now request that the codeword be transmitted again. When used for error checking, the leftmost bit is called a ***parity check bit***.

By far the most common error-detecting codes used in computers are based on the addition of a parity bit. Typically, a computer stores information in m-tuples called ***words***. Common word lengths are 8, 16, and 32 bits. One bit in the word is set aside as the parity check bit, and is not used to store information. This bit is set to either 0 or 1, depending on the number of 1s in the word.

Adding a parity check bit allows the detection of all single errors because changing a single bit either increases or decreases the number of 1s by one, and in either case the parity has been changed from even to odd, so the new word is not a codeword. (We could also construct an error detection scheme based on ***odd parity***; that is, we could set the parity check bit so that a codeword always has an odd number of 1s.)

The even parity system is easy to implement, but has two drawbacks. First, multiple errors are not detectable. Suppose an A is sent and the first and seventh bits are changed from 0 to 1. The received word is a codeword, but will be decoded into a C instead of an A. Second, we do not have the ability to correct errors. If the 8-tuple (1001 1000) is received, we know that an error has occurred, but we have no idea which bit has been changed. We will now investigate a coding scheme that will not only allow us to detect transmission errors but will actually correct the errors.

| Transmitted | Received Word | | | | | | | |
Codeword	000	001	010	011	100	101	110	111
000	0	1	1	2	1	2	2	3
111	3	2	2	1	2	1	1	0

Table 8.4: A repetition code

Example 8.5. Suppose that our original message is either a 0 or a 1, and that 0 encodes to (000) and 1 encodes to (111). If only a single error occurs during transmission, we can detect and correct the error. For example, if a (101) is received, then the second bit must have been changed from a 1 to a 0. The originally transmitted codeword must have been (111). This method will detect and correct all single errors.

In Table 8.4, we present all possible words that might be received for the transmitted codewords (000) and (111). Table 8.4 also shows the number of bits by which each received 3-tuple differs from each original codeword.

Maximum-Likelihood Decoding

The coding scheme presented in Example 8.5 is not a complete solution to the problem because it does not account for the possibility of multiple errors. For example, either a (000) or a (111) could be sent and a (001) received. We have no means of deciding from the received word whether there was a single error in the third bit or two errors, one in the first bit and one in the second. No matter what coding scheme is used, an incorrect message could be received. We could transmit a (000), have errors in all three bits, and receive the codeword (111). It is important to make explicit assumptions about the likelihood and

distribution of transmission errors so that, in a particular application, it will be known whether a given error detection scheme is appropriate. We will assume that transmission errors are rare, and, that when they do occur, they occur independently in each bit; that is, if p is the probability of an error in one bit and q is the probability of an error in a different bit, then the probability of errors occurring in both of these bits at the same time is pq. We will also assume that a received n-tuple is decoded into a codeword that is closest to it; that is, we assume that the receiver uses **maximum-likelihood decoding**.[2]
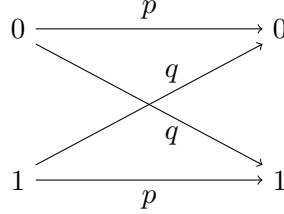
$$
\begin{array}{ccc}
0 & \xrightarrow{\quad p \quad} & 0 \\
 & q & \\
 & q & \\
1 & \xrightarrow{\quad p \quad} & 1
\end{array}
$$

Figure 8.6: Binary symmetric channel

A **binary symmetric channel** is a model that consists of a transmitter capable of sending a binary signal, either a 0 or a 1, together with a receiver. Let p be the probability that the signal is correctly received. Then $q = 1 - p$ is the probability of an incorrect reception. If a 1 is sent, then the probability that a 1 is received is p and the probability that a 0 is received is q (Figure 8.6). The probability that no errors occur during the transmission of a binary codeword of length n is p^n. For example, if $p = 0.999$ and a message consisting of 10,000 bits is sent, then the probability of a perfect transmission is

$$(0.999)^{10,000} \approx 0.00005.$$

Theorem 8.7. *If a binary n-tuple $(x_1, \ldots, x_n)$ is transmitted across a binary symmetric channel with probability p that no error will occur in each coordinate, then the probability that there are errors in exactly k coordinates is*

$$\binom{n}{k} q^k p^{n-k}.$$

PROOF. Fix k different coordinates. We first compute the probability that an error has occurred in this fixed set of coordinates. The probability of an error occurring in a particular one of these k coordinates is q; the probability that an error will not occur in any of the remaining $n - k$ coordinates is p. The probability of each of these n independent events is $q^k p^{n-k}$. The number of possible error patterns with exactly k errors occurring is equal to

$$\binom{n}{k} = \frac{n!}{k!(n-k)!},$$

the number of combinations of n things taken k at a time. Each of these error patterns has probability $q^k p^{n-k}$ of occurring; hence, the probability of all of these error patterns is

$$\binom{n}{k} q^k p^{n-k}. \qquad \square$$

---

[2]This section requires a knowledge of probability, but can be skipped without loss of continuity.

Example 8.8. Suppose that $p = 0.995$ and a 500-bit message is sent. The probability that the message was sent error-free is

$$p^n = (0.995)^{500} \approx 0.082.$$

The probability of exactly one error occurring is

$$\binom{n}{1} q p^{n-1} = 500(0.005)(0.995)^{499} \approx 0.204.$$

The probability of exactly two errors is

$$\binom{n}{2} q^2 p^{n-2} = \frac{500 \cdot 499}{2}(0.005)^2(0.995)^{498} \approx 0.257.$$

The probability of more than two errors is approximately

$$1 - 0.082 - 0.204 - 0.257 = 0.457.$$

Block Codes

If we are to develop efficient error-detecting and error-correcting codes, we will need more sophisticated mathematical tools. Group theory will allow faster methods of encoding and decoding messages. A code is an (n, m)-**block code** if the information that is to be coded can be divided into blocks of m binary digits, each of which can be encoded into n binary digits. More specifically, an (n, m)-block code consists of an **encoding function**

$$E : \mathbb{Z}_2^m \to \mathbb{Z}_2^n$$

and a **decoding function**

$$D : \mathbb{Z}_2^n \to \mathbb{Z}_2^m.$$

A **codeword** is any element in the image of E. We also require that E be one-to-one so that two information blocks will not be encoded into the same codeword. If our code is to be error-correcting, then D must be onto.

Example 8.9. The even-parity coding system developed to detect single errors in ASCII characters is an $(8, 7)$-block code. The encoding function is

$$E(x_7, x_6, \ldots, x_1) = (x_8, x_7, \ldots, x_1),$$

where $x_8 = x_7 + x_6 + \cdots + x_1$ with addition in $\mathbb{Z}_2$.

Let $\mathbf{x} = (x_1, \ldots, x_n)$ and $\mathbf{y} = (y_1, \ldots, y_n)$ be binary n-tuples. The **Hamming distance** or **distance**, $d(\mathbf{x}, \mathbf{y})$, between $\mathbf{x}$ and $\mathbf{y}$ is the number of bits in which $\mathbf{x}$ and $\mathbf{y}$ differ. The distance between two codewords is the minimum number of transmission errors required to change one codeword into the other. The **minimum distance** for a code, $d_{\min}$, is the minimum of all distances $d(\mathbf{x}, \mathbf{y})$, where $\mathbf{x}$ and $\mathbf{y}$ are distinct codewords. The **weight**, $w(\mathbf{x})$, of a binary codeword $\mathbf{x}$ is the number of 1s in $\mathbf{x}$. Clearly, $w(\mathbf{x}) = d(\mathbf{x}, \mathbf{0})$, where $\mathbf{0} = (00 \cdots 0)$.

Example 8.10. Let $\mathbf{x} = (10101)$, $\mathbf{y} = (11010)$, and $\mathbf{z} = (00011)$ be all of the codewords in some code C. Then we have the following Hamming distances:

$$d(\mathbf{x}, \mathbf{y}) = 4, \qquad d(\mathbf{x}, \mathbf{z}) = 3, \qquad d(\mathbf{y}, \mathbf{z}) = 3.$$

The minimum distance for this code is 3. We also have the following weights:

$$w(\mathbf{x}) = 3, \qquad w(\mathbf{y}) = 3, \qquad w(\mathbf{z}) = 2.$$

The following proposition lists some basic properties about the weight of a codeword and the distance between two codewords. The proof is left as an exercise.

Proposition 8.11. *Let* $\mathbf{x}$, $\mathbf{y}$, *and* $\mathbf{z}$ *be binary n-tuples. Then*

1. $w(\mathbf{x}) = d(\mathbf{x}, \mathbf{0})$;

2. $d(\mathbf{x}, \mathbf{y}) \geq 0$;

3. $d(\mathbf{x}, \mathbf{y}) = 0$ *exactly when* $\mathbf{x} = \mathbf{y}$;

4. $d(\mathbf{x}, \mathbf{y}) = d(\mathbf{y}, \mathbf{x})$;

5. $d(\mathbf{x}, \mathbf{y}) \leq d(\mathbf{x}, \mathbf{z}) + d(\mathbf{z}, \mathbf{y})$.

The weights in a particular code are usually much easier to compute than the Hamming distances between all codewords in the code. If a code is set up carefully, we can use this fact to our advantage.

Suppose that $\mathbf{x} = (1101)$ and $\mathbf{y} = (1100)$ are codewords in some code. If we transmit (1101) and an error occurs in the rightmost bit, then (1100) will be received. Since (1100) is a codeword, the decoder will decode (1100) as the transmitted message. This code is clearly not very appropriate for error detection. The problem is that $d(\mathbf{x}, \mathbf{y}) = 1$. If $\mathbf{x} = (1100)$ and $\mathbf{y} = (1010)$ are codewords, then $d(\mathbf{x}, \mathbf{y}) = 2$. If $\mathbf{x}$ is transmitted and a single error occurs, then $\mathbf{y}$ can never be received. Table 8.12 gives the distances between all 4-bit codewords in which the first three bits carry information and the fourth is an even parity check bit. We can see that the minimum distance here is 2; hence, the code is suitable as a single error-detecting code.

	0000	0011	0101	0110	1001	1010	1100	1111
0000	0	2	2	2	2	2	2	4
0011	2	0	2	2	2	2	4	2
0101	2	2	0	2	2	4	2	2
0110	2	2	2	0	4	2	2	2
1001	2	2	2	4	0	2	2	2
1010	2	2	4	2	2	0	2	2
1100	2	4	2	2	2	2	0	2
1111	4	2	2	2	2	2	2	0

Table 8.12: Distances between 4-bit codewords

To determine exactly what the error-detecting and error-correcting capabilities for a code are, we need to analyze the minimum distance for the code. Let $\mathbf{x}$ and $\mathbf{y}$ be codewords. If $d(\mathbf{x}, \mathbf{y}) = 1$ and an error occurs where $\mathbf{x}$ and $\mathbf{y}$ differ, then $\mathbf{x}$ is changed to $\mathbf{y}$. The received codeword is $\mathbf{y}$ and no error message is given. Now suppose $d(\mathbf{x}, \mathbf{y}) = 2$. Then a single error cannot change $\mathbf{x}$ to $\mathbf{y}$. Therefore, if $d_{\min} = 2$, we have the ability to detect single errors. However, suppose that $d(\mathbf{x}, \mathbf{y}) = 2$, $\mathbf{y}$ is sent, and a noncodeword $\mathbf{z}$ is received such that

$$d(\mathbf{x}, \mathbf{z}) = d(\mathbf{y}, \mathbf{z}) = 1.$$

Then the decoder cannot decide between $\mathbf{x}$ and $\mathbf{y}$. Even though we are aware that an error has occurred, we do not know what the error is.

Suppose $d_{\min} \geq 3$. Then the maximum-likelihood decoding scheme corrects all single errors. Starting with a codeword $\mathbf{x}$, an error in the transmission of a single bit gives $\mathbf{y}$

with $d(\mathbf{x}, \mathbf{y}) = 1$, but $d(\mathbf{z}, \mathbf{y}) \geq 2$ for any other codeword $\mathbf{z} \neq \mathbf{x}$. If we do not require the correction of errors, then we can detect multiple errors when a code has a minimum distance that is greater than or equal to 3.

Theorem 8.13. *Let C be a code with $d_{\min} = 2n + 1$. Then C can correct any n or fewer errors. Furthermore, any $2n$ or fewer errors can be detected in C.*

PROOF. Suppose that a codeword $\mathbf{x}$ is sent and the word $\mathbf{y}$ is received with at most n errors. Then $d(\mathbf{x}, \mathbf{y}) \leq n$. If $\mathbf{z}$ is any codeword other than $\mathbf{x}$, then

$$2n + 1 \leq d(\mathbf{x}, \mathbf{z}) \leq d(\mathbf{x}, \mathbf{y}) + d(\mathbf{y}, \mathbf{z}) \leq n + d(\mathbf{y}, \mathbf{z}).$$

Hence, $d(\mathbf{y}, \mathbf{z}) \geq n + 1$ and $\mathbf{y}$ will be correctly decoded as $\mathbf{x}$. Now suppose that $\mathbf{x}$ is transmitted and $\mathbf{y}$ is received and that at least one error has occurred, but not more than $2n$ errors. Then $1 \leq d(\mathbf{x}, \mathbf{y}) \leq 2n$. Since the minimum distance between codewords is $2n+1$, $\mathbf{y}$ cannot be a codeword. Consequently, the code can detect between 1 and $2n$ errors. $\quad\square$

Example 8.14. In Table 8.15, the codewords $\mathbf{c}_1 = (00000)$, $\mathbf{c}_2 = (00111)$, $\mathbf{c}_3 = (11100)$, and $\mathbf{c}_4 = (11011)$ determine a single error-correcting code.

	00000	00111	11100	11011
00000	0	3	3	4
00111	3	0	4	3
11100	3	4	0	3
11011	4	3	3	0

Table 8.15: Hamming distances for an error-correcting code

Historical Note

Modern coding theory began in 1948 with C. Shannon's paper, "A Mathematical Theory of Information" [7]. This paper offered an example of an algebraic code, and Shannon's Theorem proclaimed exactly how good codes could be expected to be. Richard Hamming began working with linear codes at Bell Labs in the late 1940s and early 1950s after becoming frustrated because the programs that he was running could not recover from simple errors generated by noise. Coding theory has grown tremendously in the past several decades. *The Theory of Error-Correcting Codes*, by MacWilliams and Sloane [5], published in 1977, already contained over 1500 references. Linear codes (Reed-Muller $(32, 6)$-block codes) were used on NASA's Mariner space probes. More recent space probes such as Voyager have used what are called convolution codes. Currently, very active research is being done with Goppa codes, which are heavily dependent on algebraic geometry.

8.2 Linear Codes

To gain more knowledge of a particular code and develop more efficient techniques of encoding, decoding, and error detection, we need to add additional structure to our codes. One way to accomplish this is to require that the code also be a group. A ***group code*** is a code that is also a subgroup of $\mathbb{Z}_2^n$.

To check that a code is a group code, we need only verify one thing. If we add any two elements in the code, the result must be an n-tuple that is again in the code. It is not

necessary to check that the inverse of the n-tuple is in the code, since every codeword is its own inverse, nor is it necessary to check that $\mathbf{0}$ is a codeword. For instance,

$$(11000101) + (11000101) = (00000000).$$

Example 8.16. Suppose that we have a code that consists of the following 7-tuples:

(0000000)	(0001111)	(0010101)	(0011010)
(0100110)	(0101001)	(0110011)	(0111100)
(1000011)	(1001100)	(1010110)	(1011001)
(1100101)	(1101010)	(1110000)	(1111111).

It is a straightforward though tedious task to verify that this code is also a subgroup of $\mathbb{Z}_2^7$ and, therefore, a group code. This code is a single error-detecting and single error-correcting code, but it is a long and tedious process to compute all of the distances between pairs of codewords to determine that $d_{\min} = 3$. It is much easier to see that the minimum weight of all the nonzero codewords is 3. As we will soon see, this is no coincidence. However, the relationship between weights and distances in a particular code is heavily dependent on the fact that the code is a group.

Lemma 8.17. *Let* $\mathbf{x}$ *and* $\mathbf{y}$ *be binary n-tuples. Then* $w(\mathbf{x} + \mathbf{y}) = d(\mathbf{x}, \mathbf{y})$.

PROOF.   Suppose that $\mathbf{x}$ and $\mathbf{y}$ are binary n-tuples. Then the distance between $\mathbf{x}$ and $\mathbf{y}$ is exactly the number of places in which $\mathbf{x}$ and $\mathbf{y}$ differ. But $\mathbf{x}$ and $\mathbf{y}$ differ in a particular coordinate exactly when the sum in the coordinate is 1, since

$$1 + 1 = 0$$
$$0 + 0 = 0$$
$$1 + 0 = 1$$
$$0 + 1 = 1.$$

Consequently, the weight of the sum must be the distance between the two codewords. $\square$

Theorem 8.18. *Let* $d_{\min}$ *be the minimum distance for a group code C. Then $d_{\min}$ is the minimum of all the nonzero weights of the nonzero codewords in C. That is,*

$$d_{\min} = \min\{w(\mathbf{x}) : \mathbf{x} \neq \mathbf{0}\}.$$

PROOF.   Observe that

$$\begin{aligned} d_{\min} &= \min\{d(\mathbf{x}, \mathbf{y}) : \mathbf{x} \neq \mathbf{y}\} \\ &= \min\{d(\mathbf{x}, \mathbf{y}) : \mathbf{x} + \mathbf{y} \neq \mathbf{0}\} \\ &= \min\{w(\mathbf{x} + \mathbf{y}) : \mathbf{x} + \mathbf{y} \neq \mathbf{0}\} \\ &= \min\{w(\mathbf{z}) : \mathbf{z} \neq \mathbf{0}\}. \end{aligned}$$
$\square$

Linear Codes

From Example 8.16, it is now easy to check that the minimum nonzero weight is 3; hence, the code does indeed detect and correct all single errors. We have now reduced the problem of finding "good" codes to that of generating group codes. One easy way to generate group codes is to employ a bit of matrix theory.

Define the ***inner product*** of two binary n-tuples to be

$$\mathbf{x} \cdot \mathbf{y} = x_1 y_1 + \cdots + x_n y_n,$$

where $\mathbf{x} = (x_1, x_2, \ldots, x_n)^{\mathrm{t}}$ and $\mathbf{y} = (y_1, y_2, \ldots, y_n)^{\mathrm{t}}$ are column vectors.[3] For example, if $\mathbf{x} = (011001)^{\mathrm{t}}$ and $\mathbf{y} = (110101)^{\mathrm{t}}$, then $\mathbf{x} \cdot \mathbf{y} = 0$. We can also look at an inner product as the product of a row matrix with a column matrix; that is,

$$\mathbf{x} \cdot \mathbf{y} = \mathbf{x}^{\mathrm{t}} \mathbf{y}$$

$$= \begin{pmatrix} x_1 & x_2 & \cdots & x_n \end{pmatrix} \begin{pmatrix} y_1 \\ y_2 \\ \vdots \\ y_n \end{pmatrix}$$

$$= x_1 y_1 + x_2 y_2 + \cdots + x_n y_n.$$

Example 8.19. Suppose that the words to be encoded consist of all binary 3-tuples and that our encoding scheme is even-parity. To encode an arbitrary 3-tuple, we add a fourth bit to obtain an even number of 1s. Notice that an arbitrary n-tuple $\mathbf{x} = (x_1, x_2, \ldots, x_n)^{\mathrm{t}}$ has an even number of 1s exactly when $x_1 + x_2 + \cdots + x_n = 0$; hence, a 4-tuple $\mathbf{x} = (x_1, x_2, x_3, x_4)^{\mathrm{t}}$ has an even number of 1s if $x_1 + x_2 + x_3 + x_4 = 0$, or

$$\mathbf{x} \cdot \mathbf{1} = \mathbf{x}^{\mathrm{t}} \mathbf{1} = \begin{pmatrix} x_1 & x_2 & x_3 & x_4 \end{pmatrix} \begin{pmatrix} 1 \\ 1 \\ 1 \\ 1 \end{pmatrix} = 0.$$

This example leads us to hope that there is a connection between matrices and coding theory.

Let $\mathbb{M}_{m \times n}(\mathbb{Z}_2)$ denote the set of all $m \times n$ matrices with entries in $\mathbb{Z}_2$. We do matrix operations as usual except that all our addition and multiplication operations occur in $\mathbb{Z}_2$. Define the ***null space*** of a matrix $H \in \mathbb{M}_{m \times n}(\mathbb{Z}_2)$ to be the set of all binary n-tuples $\mathbf{x}$ such that $H\mathbf{x} = \mathbf{0}$. We denote the null space of a matrix H by $\mathrm{Null}(H)$.

Example 8.20. Suppose that

$$H = \begin{pmatrix} 0 & 1 & 0 & 1 & 0 \\ 1 & 1 & 1 & 1 & 0 \\ 0 & 0 & 1 & 1 & 1 \end{pmatrix}.$$

For a 5-tuple $\mathbf{x} = (x_1, x_2, x_3, x_4, x_5)^{\mathrm{t}}$ to be in the null space of H, $H\mathbf{x} = \mathbf{0}$. Equivalently, the following system of equations must be satisfied:

$$x_2 + x_4 = 0$$
$$x_1 + x_2 + x_3 + x_4 = 0$$
$$x_3 + x_4 + x_5 = 0.$$

The set of binary 5-tuples satisfying these equations is

$$(00000) \qquad (11110) \qquad (10101) \qquad (01011).$$

This code is easily determined to be a group code.

Theorem 8.21. *Let H be in $\mathbb{M}_{m \times n}(\mathbb{Z}_2)$. Then the null space of H is a group code.*

---

[3]Since we will be working with matrices, we will write binary n-tuples as column vectors for the remainder of this chapter.

PROOF. Since each element of $\mathbb{Z}_2^n$ is its own inverse, the only thing that really needs to be checked here is closure. Let $\mathbf{x}, \mathbf{y} \in \mathrm{Null}(H)$ for some matrix H in $\mathbb{M}_{m \times n}(\mathbb{Z}_2)$. Then $H\mathbf{x} = \mathbf{0}$ and $H\mathbf{y} = \mathbf{0}$. So

$$H(\mathbf{x} + \mathbf{y}) = H\mathbf{x} + H\mathbf{y} = \mathbf{0} + \mathbf{0} = \mathbf{0}.$$

Hence, $\mathbf{x} + \mathbf{y}$ is in the null space of H and therefore must be a codeword. $\square$

A code is a ***linear code*** if it is determined by the null space of some matrix $H \in \mathbb{M}_{m \times n}(\mathbb{Z}_2)$.

Example 8.22. Let C be the code given by the matrix

$$H = \begin{pmatrix} 0 & 0 & 0 & 1 & 1 & 1 \\ 0 & 1 & 1 & 0 & 1 & 1 \\ 1 & 0 & 1 & 0 & 0 & 1 \end{pmatrix}.$$

Suppose that the 6-tuple $\mathbf{x} = (010011)^{\mathrm{t}}$ is received. It is a simple matter of matrix multiplication to determine whether or not $\mathbf{x}$ is a codeword. Since

$$H\mathbf{x} = \begin{pmatrix} 0 \\ 1 \\ 1 \end{pmatrix},$$

the received word is not a codeword. We must either attempt to correct the word or request that it be transmitted again.

8.3 Parity-Check and Generator Matrices

We need to find a systematic way of generating linear codes as well as fast methods of decoding. By examining the properties of a matrix H and by carefully choosing H, it is possible to develop very efficient methods of encoding and decoding messages. To this end, we will introduce standard generator and canonical parity-check matrices.

Suppose that H is an $m \times n$ matrix with entries in $\mathbb{Z}_2$ and $n > m$. If the last m columns of the matrix form the $m \times m$ identity matrix, I_m, then the matrix is a ***canonical parity-check matrix***. More specifically, $H = (A \mid I_m)$, where A is the $m \times (n - m)$ matrix

$$\begin{pmatrix} a_{11} & a_{12} & \cdots & a_{1,n-m} \\ a_{21} & a_{22} & \cdots & a_{2,n-m} \\ \vdots & \vdots & \ddots & \vdots \\ a_{m1} & a_{m2} & \cdots & a_{m,n-m} \end{pmatrix}$$

and I_m is the $m \times m$ identity matrix

$$\begin{pmatrix} 1 & 0 & \cdots & 0 \\ 0 & 1 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & 1 \end{pmatrix}.$$

With each canonical parity-check matrix we can associate an $n \times (n - m)$ ***standard generator matrix***

$$G = \left(\frac{I_{n-m}}{A} \right).$$

Our goal will be to show that an $\mathbf{x}$ satisfying $G\mathbf{x} = \mathbf{y}$ exists if and only if $H\mathbf{y} = \mathbf{0}$. Given a message block $\mathbf{x}$ to be encoded, the matrix G will allow us to quickly encode it into a linear codeword $\mathbf{y}$.

Example 8.23. Suppose that we have the following eight words to be encoded:

$$(000), (001), (010), \ldots, (111).$$

For

$$A = \begin{pmatrix} 0 & 1 & 1 \\ 1 & 1 & 0 \\ 1 & 0 & 1 \end{pmatrix},$$

the associated standard generator and canonical parity-check matrices are

$$G = \begin{pmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1 \\ 0 & 1 & 1 \\ 1 & 1 & 0 \\ 1 & 0 & 1 \end{pmatrix}$$

and

$$H = \begin{pmatrix} 0 & 1 & 1 & 1 & 0 & 0 \\ 1 & 1 & 0 & 0 & 1 & 0 \\ 1 & 0 & 1 & 0 & 0 & 1 \end{pmatrix},$$

respectively.

Observe that the rows in H represent the parity checks on certain bit positions in a 6-tuple. The 1s in the identity matrix serve as parity checks for the 1s in the same row. If $\mathbf{x} = (x_1, x_2, x_3, x_4, x_5, x_6)$, then

$$\mathbf{0} = H\mathbf{x} = \begin{pmatrix} x_2 + x_3 + x_4 \\ x_1 + x_2 + x_5 \\ x_1 + x_3 + x_6 \end{pmatrix},$$

which yields a system of equations:

$$x_2 + x_3 + x_4 = 0$$
$$x_1 + x_2 + x_5 = 0$$
$$x_1 + x_3 + x_6 = 0.$$

Here x_4 serves as a check bit for x_2 and x_3; x_5 is a check bit for x_1 and x_2; and x_6 is a check bit for x_1 and x_3. The identity matrix keeps x_4, x_5, and x_6 from having to check on each other. Hence, x_1, x_2, and x_3 can be arbitrary but x_4, x_5, and x_6 must be chosen to ensure parity. The null space of H is easily computed to be

$$(000000) \quad (001101) \quad (010110) \quad (011011)$$
$$(100011) \quad (101110) \quad (110101) \quad (111000).$$

An even easier way to compute the null space is with the generator matrix G (Table 8.24).

Message Word $\mathbf{x}$	Codeword $G\mathbf{x}$
000	000000
001	001101
010	010110
011	011011
100	100011
101	101110
110	110101
111	111000

Table 8.24: A matrix-generated code

Theorem 8.25. *If $H \in \mathbb{M}_{m \times n}(\mathbb{Z}_2)$ is a canonical parity-check matrix, then* $\mathrm{Null}(H)$ *consists of all $\mathbf{x} \in \mathbb{Z}_2^n$ whose first $n - m$ bits are arbitrary but whose last m bits are determined by $H\mathbf{x} = \mathbf{0}$. Each of the last m bits serves as an even parity check bit for some of the first $n - m$ bits. Hence, H gives rise to an $(n, n - m)$-block code.*

We leave the proof of this theorem as an exercise. In light of the theorem, the first $n - m$ bits in $\mathbf{x}$ are called ***information bits*** and the last m bits are called ***check bits***. In Example 8.23, the first three bits are the information bits and the last three are the check bits.

Theorem 8.26. *Suppose that G is an $n \times k$ standard generator matrix. Then $C = \left\{ \mathbf{y} : G\mathbf{x} = \mathbf{y} \text{ for } \mathbf{x} \in \mathbb{Z}_2^k \right\}$ is an (n, k)-block code. More specifically, C is a group code.*

PROOF.   Let $G\mathbf{x}_1 = \mathbf{y}_1$ and $G\mathbf{x}_2 = \mathbf{y}_2$ be two codewords. Then $\mathbf{y}_1 + \mathbf{y}_2$ is in C since

$$G(\mathbf{x}_1 + \mathbf{x}_2) = G\mathbf{x}_1 + G\mathbf{x}_2 = \mathbf{y}_1 + \mathbf{y}_2.$$

We must also show that two message blocks cannot be encoded into the same codeword. That is, we must show that if $G\mathbf{x} = G\mathbf{y}$, then $\mathbf{x} = \mathbf{y}$. Suppose that $G\mathbf{x} = G\mathbf{y}$. Then

$$G\mathbf{x} - G\mathbf{y} = G(\mathbf{x} - \mathbf{y}) = \mathbf{0}.$$

However, the first k coordinates in $G(\mathbf{x} - \mathbf{y})$ are exactly $x_1 - y_1, \ldots, x_k - y_k$, since they are determined by the identity matrix, I_k, part of G. Hence, $G(\mathbf{x} - \mathbf{y}) = \mathbf{0}$ exactly when $\mathbf{x} = \mathbf{y}$. $\qquad\square$

Before we can prove the relationship between canonical parity-check matrices and standard generating matrices, we need to prove a lemma.

Lemma 8.27. *Let $H = (A \mid I_m)$ be an $m \times n$ canonical parity-check matrix and $G = \left(\frac{I_{n-m}}{A} \right)$ be the corresponding $n \times (n - m)$ standard generator matrix. Then $HG = \mathbf{0}$.*

PROOF.   Let $C = HG$. The ijth entry in C is

$$
\begin{aligned}
c_{ij} &= \sum_{k=1}^{n} h_{ik} g_{kj} \\
&= \sum_{k=1}^{n-m} h_{ik} g_{kj} + \sum_{k=n-m+1}^{n} h_{ik} g_{kj} \\
&= \sum_{k=1}^{n-m} a_{ik} \delta_{kj} + \sum_{k=n-m+1}^{n} \delta_{i-(m-n),k} a_{kj}
\end{aligned}
$$

$$= a_{ij} + a_{ij}$$
$$= 0,$$

where

$$\delta_{ij} = \begin{cases} 1, & i = j \\ 0, & i \neq j \end{cases}$$

is the Kronecker delta. $\square$

Theorem 8.28. *Let $H = (A \mid I_m)$ be an $m \times n$ canonical parity-check matrix and let $G = \left(\frac{I_{n-m}}{A} \right)$ be the $n \times (n-m)$ standard generator matrix associated with H. Let C be the code generated by G. Then $\mathbf{y}$ is in C if and only if $H\mathbf{y} = \mathbf{0}$. In particular, C is a linear code with canonical parity-check matrix H.*

PROOF. First suppose that $\mathbf{y} \in C$. Then $G\mathbf{x} = \mathbf{y}$ for some $\mathbf{x} \in \mathbb{Z}_2^m$. By Lemma 8.27, $H\mathbf{y} = HG\mathbf{x} = \mathbf{0}$.

Conversely, suppose that $\mathbf{y} = (y_1, \ldots, y_n)^{\mathrm{t}}$ is in the null space of H. We need to find an $\mathbf{x}$ in $\mathbb{Z}_2^{n-m}$ such that $G\mathbf{x}^{\mathrm{t}} = \mathbf{y}$. Since $H\mathbf{y} = \mathbf{0}$, the following set of equations must be satisfied:

$$a_{11}y_1 + a_{12}y_2 + \cdots + a_{1,n-m}y_{n-m} + y_{n-m+1} = 0$$
$$a_{21}y_1 + a_{22}y_2 + \cdots + a_{2,n-m}y_{n-m} + y_{n-m+1} = 0$$
$$\vdots$$
$$a_{m1}y_1 + a_{m2}y_2 + \cdots + a_{m,n-m}y_{n-m} + y_{n-m+1} = 0.$$

Equivalently, $y_{n-m+1}, \ldots, y_n$ are determined by $y_1, \ldots, y_{n-m}$:

$$y_{n-m+1} = a_{11}y_1 + a_{12}y_2 + \cdots + a_{1,n-m}y_{n-m}$$
$$y_{n-m+1} = a_{21}y_1 + a_{22}y_2 + \cdots + a_{2,n-m}y_{n-m}$$
$$\vdots$$
$$y_{n-m+1} = a_{m1}y_1 + a_{m2}y_2 + \cdots + a_{m,n-m}y_{n-m}.$$

Consequently, we can let $x_i = y_i$ for $i = 1, \ldots, n - m$. $\square$

It would be helpful if we could compute the minimum distance of a linear code directly from its matrix H in order to determine the error-detecting and error-correcting capabilities of the code. Suppose that

$$\mathbf{e}_1 = (100 \cdots 00)^{\mathrm{t}}$$
$$\mathbf{e}_2 = (010 \cdots 00)^{\mathrm{t}}$$
$$\vdots$$
$$\mathbf{e}_n = (000 \cdots 01)^{\mathrm{t}}$$

are the n-tuples in $\mathbb{Z}_2^n$ of weight 1. For an $m \times n$ binary matrix H, $H\mathbf{e}_i$ is exactly the ith column of the matrix H.

Example 8.29. Observe that

$$\begin{pmatrix} 1 & 1 & 1 & 0 & 0 \\ 1 & 0 & 0 & 1 & 0 \\ 1 & 1 & 0 & 0 & 1 \end{pmatrix} \begin{pmatrix} 0 \\ 1 \\ 0 \\ 0 \\ 0 \end{pmatrix} = \begin{pmatrix} 1 \\ 0 \\ 1 \end{pmatrix}.$$

We state this result in the following proposition and leave the proof as an exercise.

Proposition 8.30. *Let $\mathbf{e}_i$ be the binary n-tuple with a 1 in the ith coordinate and 0's elsewhere and suppose that $H \in \mathbb{M}_{m \times n}(\mathbb{Z}_2)$. Then $H\mathbf{e}_i$ is the ith column of the matrix H.*

Theorem 8.31. *Let H be an $m \times n$ binary matrix. Then the null space of H is a single error-detecting code if and only if no column of H consists entirely of zeros.*

PROOF.   Suppose that $\mathrm{Null}(H)$ is a single error-detecting code. Then the minimum distance of the code must be at least 2. Since the null space is a group code, it is sufficient to require that the code contain no codewords of less than weight 2 other than the zero codeword. That is, $\mathbf{e}_i$ must not be a codeword for $i = 1, \ldots, n$. Since $H\mathbf{e}_i$ is the ith column of H, the only way in which $\mathbf{e}_i$ could be in the null space of H would be if the ith column were all zeros, which is impossible; hence, the code must have the capability to detect at least single errors.

Conversely, suppose that no column of H is the zero column. By Proposition 8.30, $H\mathbf{e}_i \neq \mathbf{0}$. $\qquad\square$

Example 8.32. If we consider the matrices

$$H_1 = \begin{pmatrix} 1 & 1 & 1 & 0 & 0 \\ 1 & 0 & 0 & 1 & 0 \\ 1 & 1 & 0 & 0 & 1 \end{pmatrix}$$

and

$$H_2 = \begin{pmatrix} 1 & 1 & 1 & 0 & 0 \\ 1 & 0 & 0 & 0 & 0 \\ 1 & 1 & 0 & 0 & 1 \end{pmatrix},$$

then the null space of H_1 is a single error-detecting code and the null space of H_2 is not.

We can even do better than Theorem 8.31. This theorem gives us conditions on a matrix H that tell us when the minimum weight of the code formed by the null space of H is 2. We can also determine when the minimum distance of a linear code is 3 by examining the corresponding matrix.

Example 8.33. If we let

$$H = \begin{pmatrix} 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 \\ 1 & 1 & 0 & 0 \end{pmatrix}$$

and want to determine whether or not H is the canonical parity-check matrix for an error-correcting code, it is necessary to make certain that $\mathrm{Null}(H)$ does not contain any 4-tuples of weight 2. That is, (1100), (1010), (1001), (0110), (0101), and (0011) must not be in $\mathrm{Null}(H)$. The next theorem states that we can indeed determine that the code generated by H is error-correcting by examining the columns of H. Notice in this example that not only does H have no zero columns, but also that no two columns are the same.

Theorem 8.34. *Let H be a binary matrix. The null space of H is a single error-correcting code if and only if H does not contain any zero columns and no two columns of H are identical.*

PROOF.   The n-tuple $\mathbf{e}_i + \mathbf{e}_j$ has 1s in the ith and jth entries and 0s elsewhere, and $w(\mathbf{e}_i + \mathbf{e}_j) = 2$ for $i \neq j$. Since

$$\mathbf{0} = H(\mathbf{e}_i + \mathbf{e}_j) = H\mathbf{e}_i + H\mathbf{e}_j$$

can only occur if the ith and jth columns are identical, the null space of H is a single error-correcting code. $\qquad\square$

Suppose now that we have a canonical parity-check matrix H with three rows. Then we might ask how many more columns we can add to the matrix and still have a null space that is a single error-detecting and single error-correcting code. Since each column has three entries, there are $2^3 = 8$ possible distinct columns. We cannot add the columns

$$\begin{pmatrix} 0 \\ 0 \\ 0 \end{pmatrix}, \begin{pmatrix} 1 \\ 0 \\ 0 \end{pmatrix}, \begin{pmatrix} 0 \\ 1 \\ 0 \end{pmatrix}, \begin{pmatrix} 0 \\ 0 \\ 1 \end{pmatrix}.$$

So we can add as many as four columns and still maintain a minimum distance of 3.

In general, if H is an $m \times n$ canonical parity-check matrix, then there are $n - m$ information positions in each codeword. Each column has m bits, so there are 2^m possible distinct columns. It is necessary that the columns $\mathbf{0}, \mathbf{e}_1, \ldots, \mathbf{e}_m$ be excluded, leaving $2^m - (1 + m)$ remaining columns for information if we are still to maintain the ability not only to detect but also to correct single errors.

8.4 Efficient Decoding

We are now at the stage where we are able to generate linear codes that detect and correct errors fairly easily, but it is still a time-consuming process to decode a received n-tuple and determine which is the closest codeword, because the received n-tuple must be compared to each possible codeword to determine the proper decoding. This can be a serious impediment if the code is very large.

Example 8.35. Given the binary matrix

$$H = \begin{pmatrix} 1 & 1 & 1 & 0 & 0 \\ 0 & 1 & 0 & 1 & 0 \\ 1 & 0 & 0 & 0 & 1 \end{pmatrix}$$

and the 5-tuples $\mathbf{x} = (11011)^{\mathrm{t}}$ and $\mathbf{y} = (01011)^{\mathrm{t}}$, we can compute

$$H\mathbf{x} = \begin{pmatrix} 0 \\ 0 \\ 0 \end{pmatrix} \qquad \text{and} \qquad H\mathbf{y} = \begin{pmatrix} 1 \\ 0 \\ 1 \end{pmatrix}.$$

Hence, $\mathbf{x}$ is a codeword and $\mathbf{y}$ is not, since $\mathbf{x}$ is in the null space and $\mathbf{y}$ is not. Notice that $H\mathbf{y}$ is identical to the first column of H. In fact, this is where the error occurred. If we flip the first bit in $\mathbf{y}$ from 0 to 1, then we obtain $\mathbf{x}$.

If H is an $m \times n$ matrix and $\mathbf{x} \in \mathbb{Z}_2^n$, then we say that the **syndrome** of $\mathbf{x}$ is $H\mathbf{x}$. The following proposition allows the quick detection and correction of errors.

Proposition 8.36. *Let the $m \times n$ binary matrix H determine a linear code and let $\mathbf{x}$ be the received n-tuple. Write $\mathbf{x}$ as $\mathbf{x} = \mathbf{c} + \mathbf{e}$, where $\mathbf{c}$ is the transmitted codeword and $\mathbf{e}$ is the transmission error. Then the syndrome $H\mathbf{x}$ of the received codeword $\mathbf{x}$ is also the syndrome of the error $\mathbf{e}$.*

PROOF. The proof follows from the fact that

$$H\mathbf{x} = H(\mathbf{c} + \mathbf{e}) = H\mathbf{c} + H\mathbf{e} = \mathbf{0} + H\mathbf{e} = H\mathbf{e}. \qquad \square$$

This proposition tells us that the syndrome of a received word depends solely on the error and not on the transmitted codeword. The proof of the following theorem follows immediately from Proposition 8.36 and from the fact that $H\mathbf{e}$ is the ith column of the matrix H.

Theorem 8.37. *Let* $H \in \mathbb{M}_{m \times n}(\mathbb{Z}_2)$ *and suppose that the linear code corresponding to* H *is single error-correcting. Let* $\mathbf{r}$ *be a received* n-*tuple that was transmitted with at most one error. If the syndrome of* $\mathbf{r}$ *is* $\mathbf{0}$, *then no error has occurred; otherwise, if the syndrome of* $\mathbf{r}$ *is equal to some column of* H, *say the* ith *column, then the error has occurred in the* ith *bit.*

Example 8.38. Consider the matrix

$$H = \begin{pmatrix} 1 & 0 & 1 & 1 & 0 & 0 \\ 0 & 1 & 1 & 0 & 1 & 0 \\ 1 & 1 & 1 & 0 & 0 & 1 \end{pmatrix}$$

and suppose that the 6-tuples $\mathbf{x} = (111110)^{\mathrm{t}}$, $\mathbf{y} = (111111)^{\mathrm{t}}$, and $\mathbf{z} = (010111)^{\mathrm{t}}$ have been received. Then

$$H\mathbf{x} = \begin{pmatrix} 1 \\ 1 \\ 1 \end{pmatrix}, H\mathbf{y} = \begin{pmatrix} 1 \\ 1 \\ 0 \end{pmatrix}, H\mathbf{z} = \begin{pmatrix} 1 \\ 0 \\ 0 \end{pmatrix}.$$

Hence, $\mathbf{x}$ has an error in the third bit and $\mathbf{z}$ has an error in the fourth bit. The transmitted codewords for $\mathbf{x}$ and $\mathbf{z}$ must have been (110110) and (010011), respectively. The syndrome of $\mathbf{y}$ does not occur in any of the columns of the matrix H, so multiple errors must have occurred to produce $\mathbf{y}$.

Coset Decoding

We can use group theory to obtain another way of decoding messages. A linear code C is a subgroup of $\mathbb{Z}_2^n$. **Coset** or **standard decoding** uses the cosets of C in $\mathbb{Z}_2^n$ to implement maximum-likelihood decoding. Suppose that C is an (n, m)-linear code. A coset of C in $\mathbb{Z}_2^n$ is written in the form $\mathbf{x} + C$, where $\mathbf{x} \in \mathbb{Z}_2^n$. By Lagrange's Theorem (Theorem 6.10), there are 2^{n-m} distinct cosets of C in $\mathbb{Z}_2^n$.

Example 8.39. Let C be the $(5, 3)$-linear code given by the parity-check matrix

$$H = \begin{pmatrix} 0 & 1 & 1 & 0 & 0 \\ 1 & 0 & 0 & 1 & 0 \\ 1 & 1 & 0 & 0 & 1 \end{pmatrix}.$$

The code consists of the codewords

$$(00000) \quad (01101) \quad (10011) \quad (11110).$$

There are $2^{5-2} = 2^3$ cosets of C in $\mathbb{Z}_2^5$, each with order $2^2 = 4$. These cosets are listed in Table 8.40.

Coset Representative	Coset
C	$(00000)(01101)(10011)(11110)$
$(10000) + C$	$(10000)(11101)(00011)(01110)$
$(01000) + C$	$(01000)(00101)(11011)(10110)$
$(00100) + C$	$(00100)(01001)(10111)(11010)$
$(00010) + C$	$(00010)(01111)(10001)(11100)$
$(00001) + C$	$(00001)(01100)(10010)(11111)$
$(10100) + C$	$(00111)(01010)(10100)(11001)$
$(00110) + C$	$(00110)(01011)(10101)(11000)$

Table 8.40: Cosets of C

Our task is to find out how knowing the cosets might help us to decode a message. Suppose that $\mathbf{x}$ was the original codeword sent and that $\mathbf{r}$ is the n-tuple received. If $\mathbf{e}$ is the transmission error, then $\mathbf{r} = \mathbf{e} + \mathbf{x}$ or, equivalently, $\mathbf{x} = \mathbf{e} + \mathbf{r}$. However, this is exactly the statement that $\mathbf{r}$ is an element in the coset $\mathbf{e} + C$. In maximum-likelihood decoding we expect the error $\mathbf{e}$ to be as small as possible; that is, $\mathbf{e}$ will have the least weight. An n-tuple of least weight in a coset is called a ***coset leader***. Once we have determined a coset leader for each coset, the decoding process becomes a task of calculating $\mathbf{r} + \mathbf{e}$ to obtain $\mathbf{x}$.

Example 8.41. In Table 8.40, notice that we have chosen a representative of the least possible weight for each coset. These representatives are coset leaders. Now suppose that $\mathbf{r} = (01111)$ is the received word. To decode $\mathbf{r}$, we find that it is in the coset $(00010) + C$; hence, the originally transmitted codeword must have been $(01101) = (01111) + (00010)$.

A potential problem with this method of decoding is that we might have to examine every coset for the received codeword. The following proposition gives a method of implementing coset decoding. It states that we can associate a syndrome with each coset; hence, we can make a table that designates a coset leader corresponding to each syndrome. Such a list is called a ***decoding table***.

Syndrome	Coset Leader
(000)	(00000)
(001)	(00001)
(010)	(00010)
(011)	(10000)
(100)	(00100)
(101)	(01000)
(110)	(00110)
(111)	(10100)

Table 8.42: Syndromes for each coset

Proposition 8.43. *Let C be an (n, k)-linear code given by the matrix H and suppose that $\mathbf{x}$ and $\mathbf{y}$ are in $\mathbb{Z}_2^n$. Then $\mathbf{x}$ and $\mathbf{y}$ are in the same coset of C if and only if $H\mathbf{x} = H\mathbf{y}$. That is, two n-tuples are in the same coset if and only if their syndromes are the same.*

PROOF. Two n-tuples $\mathbf{x}$ and $\mathbf{y}$ are in the same coset of C exactly when $\mathbf{x} - \mathbf{y} \in C$; however, this is equivalent to $H(\mathbf{x} - \mathbf{y}) = 0$ or $H\mathbf{x} = H\mathbf{y}$. $\square$

Example 8.44. Table 8.42 is a decoding table for the code C given in Example 8.39. If $\mathbf{x} = (01111)$ is received, then its syndrome can be computed to be

$$H\mathbf{x} = \begin{pmatrix} 0 \\ 1 \\ 1 \end{pmatrix}.$$

Examining the decoding table, we determine that the coset leader is (00010). It is now easy to decode the received codeword.

Given an (n, k)-block code, the question arises of whether or not coset decoding is a manageable scheme. A decoding table requires a list of cosets and syndromes, one for each of the 2^{n-k} cosets of C. Suppose that we have a $(32, 24)$-block code. We have a huge number of codewords, 2^{24}, yet there are only $2^{32-24} = 2^8 = 256$ cosets.

8.5   Exercises

1.   Why is the following encoding scheme not acceptable?

Information	0	1	2	3	4	5	6	7	8
Codeword	000	001	010	011	101	110	111	000	001

2.   Without doing any addition, explain why the following set of 4-tuples in $\mathbb{Z}_2^4$ cannot be a group code.

$$(0110) \quad (1001) \quad (1010) \quad (1100)$$

3.   Compute the Hamming distances between the following pairs of n-tuples.

(a) $(011010), (011100)$ (c) $(00110), (01111)$

(b) $(11110101), (01010100)$ (d) $(1001), (0111)$

4.   Compute the weights of the following n-tuples.

(a) (011010) (c) (01111)

(b) (11110101) (d) (1011)

5.   Suppose that a linear code C has a minimum weight of 7. What are the error-detection and error-correction capabilities of C?

6.   In each of the following codes, what is the minimum distance for the code? What is the best situation we might hope for in connection with error detection and error correction?

(a) (011010) (011100) (110111) (110000)

(b) (011100) (011011) (111011) (100011)
      (000000) (010101) (110100) (110011)

(c) (000000) (011100) (110101) (110001)

(d) (0110110) (0111100) (1110000) (1111111)
      (1001001) (1000011) (0001111) (0000000)

7.   Compute the null space of each of the following matrices. What type of (n, k)-block codes are the null spaces? Can you find a matrix (not necessarily a standard generator matrix) that generates each code? Are your generator matrices unique?

(a)
$$\begin{pmatrix} 0 & 1 & 0 & 0 & 0 \\ 1 & 0 & 1 & 0 & 1 \\ 1 & 0 & 0 & 1 & 0 \end{pmatrix}$$

(b)
$$\begin{pmatrix} 1 & 0 & 1 & 0 & 0 & 0 \\ 1 & 1 & 0 & 1 & 0 & 0 \\ 0 & 1 & 0 & 0 & 1 & 0 \\ 1 & 1 & 0 & 0 & 0 & 1 \end{pmatrix}$$

(c)
$$\begin{pmatrix} 1 & 0 & 0 & 1 & 1 \\ 0 & 1 & 0 & 1 & 1 \end{pmatrix}$$

(d)
$$\begin{pmatrix} 0 & 0 & 0 & 1 & 1 & 1 & 1 \\ 0 & 1 & 1 & 0 & 0 & 1 & 1 \\ 1 & 0 & 1 & 0 & 1 & 0 & 1 \\ 0 & 1 & 1 & 0 & 0 & 1 & 1 \end{pmatrix}$$

8. Construct a $(5, 2)$-block code. Discuss both the error-detection and error-correction capabilities of your code.

9. Let C be the code obtained from the null space of the matrix

$$H = \begin{pmatrix} 0 & 1 & 0 & 0 & 1 \\ 1 & 0 & 1 & 0 & 1 \\ 0 & 0 & 1 & 1 & 1 \end{pmatrix}.$$

Decode the message

$$01111 \quad 10101 \quad 01110 \quad 00011$$

if possible.

10. Suppose that a 1000-bit binary message is transmitted. Assume that the probability of a single error is p and that the errors occurring in different bits are independent of one another. If $p = 0.01$, what is the probability of more than one error occurring? What is the probability of exactly two errors occurring? Repeat this problem for $p = 0.0001$.

11. Which matrices are canonical parity-check matrices? For those matrices that are canonical parity-check matrices, what are the corresponding standard generator matrices? What are the error-detection and error-correction capabilities of the code generated by each of these matrices?

(a)
$$\begin{pmatrix} 1 & 1 & 0 & 0 & 0 \\ 0 & 0 & 1 & 0 & 0 \\ 0 & 0 & 0 & 1 & 0 \\ 1 & 0 & 0 & 0 & 1 \end{pmatrix}$$

(b)
$$\begin{pmatrix} 0 & 1 & 1 & 0 & 0 & 0 \\ 1 & 1 & 0 & 1 & 0 & 0 \\ 0 & 1 & 0 & 0 & 1 & 0 \\ 1 & 1 & 0 & 0 & 0 & 1 \end{pmatrix}$$

(c)
$$\begin{pmatrix} 1 & 1 & 1 & 0 \\ 1 & 0 & 0 & 1 \end{pmatrix}$$

(d)
$$\begin{pmatrix} 0 & 0 & 0 & 1 & 0 & 0 & 0 \\ 0 & 1 & 1 & 0 & 1 & 0 & 0 \\ 1 & 0 & 1 & 0 & 0 & 1 & 0 \\ 0 & 1 & 1 & 0 & 0 & 0 & 1 \end{pmatrix}$$

12. List all possible syndromes for the codes generated by each of the matrices in Exercise 8.5.11.

13. Let

$$H = \begin{pmatrix} 0 & 1 & 1 & 1 & 1 \\ 0 & 0 & 0 & 1 & 1 \\ 1 & 0 & 1 & 0 & 1 \end{pmatrix}.$$

Compute the syndrome caused by each of the following transmission errors.

(a) An error in the first bit.

(b) An error in the third bit.

(c) An error in the last bit.

(d) Errors in the third and fourth bits.

14.   Let C be the group code in $\mathbb{Z}_2^3$ defined by the codewords (000) and (111). Compute the cosets of C in $\mathbb{Z}_2^3$. Why was there no need to specify right or left cosets? Give the single transmission error, if any, to which each coset corresponds.

15.   For each of the following matrices, find the cosets of the corresponding code C. Give a decoding table for each code if possible.

(a)

$$\begin{pmatrix} 0 & 1 & 0 & 0 & 0 \\ 1 & 0 & 1 & 0 & 1 \\ 1 & 0 & 0 & 1 & 0 \end{pmatrix}$$

(c)

$$\begin{pmatrix} 1 & 0 & 0 & 1 & 1 \\ 0 & 1 & 0 & 1 & 1 \end{pmatrix}$$

(b)

$$\begin{pmatrix} 0 & 0 & 1 & 0 & 0 \\ 1 & 1 & 0 & 1 & 0 \\ 0 & 1 & 0 & 1 & 0 \\ 1 & 1 & 0 & 0 & 1 \end{pmatrix}$$

(d)

$$\begin{pmatrix} 1 & 0 & 0 & 1 & 1 & 1 & 1 \\ 1 & 1 & 1 & 0 & 0 & 1 & 1 \\ 1 & 0 & 1 & 0 & 1 & 0 & 1 \\ 1 & 1 & 1 & 0 & 0 & 1 & 0 \end{pmatrix}$$

16.   Let $\mathbf{x}$, $\mathbf{y}$, and $\mathbf{z}$ be binary n-tuples. Prove each of the following statements.

(a) $w(\mathbf{x}) = d(\mathbf{x}, \mathbf{0})$

(b) $d(\mathbf{x}, \mathbf{y}) = d(\mathbf{x} + \mathbf{z}, \mathbf{y} + \mathbf{z})$

(c) $d(\mathbf{x}, \mathbf{y}) = w(\mathbf{x} - \mathbf{y})$

17.   A **metric** on a set X is a map $d : X \times X \to \mathbb{R}$ satisfying the following conditions.

(a) $d(\mathbf{x}, \mathbf{y}) \geq 0$ for all $\mathbf{x}, \mathbf{y} \in X$;

(b) $d(\mathbf{x}, \mathbf{y}) = 0$ exactly when $\mathbf{x} = \mathbf{y}$;

(c) $d(\mathbf{x}, \mathbf{y}) = d(\mathbf{y}, \mathbf{x})$;

(d) $d(\mathbf{x}, \mathbf{y}) \leq d(\mathbf{x}, \mathbf{z}) + d(\mathbf{z}, \mathbf{y})$.

In other words, a metric is simply a generalization of the notion of distance. Prove that Hamming distance is a metric on $\mathbb{Z}_2^n$. Decoding a message actually reduces to deciding which is the closest codeword in terms of distance.

18.   Let C be a linear code. Show that either the ith coordinates in the codewords of C are all zeros or exactly half of them are zeros.

19.   Let C be a linear code. Show that either every codeword has even weight or exactly half of the codewords have even weight.

20.   Show that the codewords of even weight in a linear code C are also a linear code.

21.   If we are to use an error-correcting linear code to transmit the 128 ASCII characters, what size matrix must be used? What size matrix must be used to transmit the extended ASCII character set of 256 characters? What if we require only error detection in both cases?

22.   Find the canonical parity-check matrix that gives the even parity check bit code with three information positions. What is the matrix for seven information positions? What are the corresponding standard generator matrices?

23.   How many check positions are needed for a single error-correcting code with 20 information positions? With 32 information positions?

24.   Let $\mathbf{e}_i$ be the binary n-tuple with a 1 in the ith coordinate and 0's elsewhere and suppose that $H \in \mathbb{M}_{m \times n}(\mathbb{Z}_2)$. Show that $H\mathbf{e}_i$ is the ith column of the matrix H.

25.   Let C be an (n, k)-linear code. Define the ***dual*** or ***orthogonal code*** of C to be

$$C^{\perp} = \{\mathbf{x} \in \mathbb{Z}_2^n : \mathbf{x} \cdot \mathbf{y} = 0 \text{ for all } \mathbf{y} \in C\}.$$

(a) Find the dual code of the linear code C where C is given by the matrix

$$\begin{pmatrix} 1 & 1 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 & 1 \\ 1 & 0 & 0 & 1 & 0 \end{pmatrix}.$$

(b) Show that $C^{\perp}$ is an $(n, n - k)$-linear code.

(c) Find the standard generator and parity-check matrices of C and $C^{\perp}$. What happens in general? Prove your conjecture.

26.   Let H be an $m \times n$ matrix over $\mathbb{Z}_2$, where the ith column is the number i written in binary with m bits. The null space of such a matrix is called a ***Hamming code***.

(a) Show that the matrix

$$H = \begin{pmatrix} 0 & 0 & 0 & 1 & 1 & 1 \\ 0 & 1 & 1 & 0 & 0 & 1 \\ 1 & 0 & 1 & 0 & 1 & 0 \end{pmatrix}$$

generates a Hamming code. What are the error-correcting properties of a Hamming code?

(b) The column corresponding to the syndrome also marks the bit that was in error; that is, the ith column of the matrix is i written as a binary number, and the syndrome immediately tells us which bit is in error. If the received word is (101011), compute the syndrome. In which bit did the error occur in this case, and what codeword was originally transmitted?

(c) Give a binary matrix H for the Hamming code with six information positions and four check positions. What are the check positions and what are the information positions? Encode the messages (101101) and (001001). Decode the received words (0010000101) and (0000101100). What are the possible syndromes for this code?

(d) What is the number of check bits and the number of information bits in an (m, n)-block Hamming code? Give both an upper and a lower bound on the number of information bits in terms of the number of check bits. Hamming codes having the maximum possible number of information bits with k check bits are called ***perfect***. Every possible syndrome except **0** occurs as a column. If the number of information bits is less than the maximum, then the code is called ***shortened***. In this case, give an example showing that some syndromes can represent multiple errors.

8.6   Programming Exercises

1.   Write a program to implement a $(16, 12)$-linear code. Your program should be able to encode and decode messages using coset decoding. Once your program is written, write a program to simulate a binary symmetric channel with transmission noise. Compare the results of your simulation with the theoretically predicted error probability.

References and Suggested Readings

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8.7   Sage

Sage has a full suite of linear codes and a variety of methods that may be used to investigate them.

Constructing Linear Codes

The `codes` object can be used to get a concise listing of the available implemented codes. Type `codes.` and press the `Tab` key and most interfaces to Sage will give you a list. You can then use a question mark at the end of a method name to learn about the various parameters.

```
codes.
```

We will use the classic binary Hamming $(7, 4)$ code as an illustration. "Binary" means we have vectors with just 0's and 1's, the 7 is the length and means the vectors have 7 coordinates, and the 4 is the dimension, meaning this code has $2^4 = 16$ vectors comprising the code. The documentation assumes we know a few things from later in the course. We use `GF(2)` to specify that our code is binary — this will make more sense at the end of the course. A second parameter is `r` and we can see from the formulas in the documenation that setting `r=3` will give length 7.

```
H = codes.HammingCode(GF(2), 3); H
```

```
[7, 4] Hamming Code over GF(2)
```

Properties of Linear Codes

We can examine the Hamming code we just built. First the dimension.

```
H.dimension()
```

```
4
```

The code is small enough that we can list all the codewords.

```
H.list()
```

```
[(0, 0, 0, 0, 0, 0, 0), (1, 0, 0, 0, 0, 1, 1), (0, 1, 0, 0, 1, 0, 1),
 (1, 1, 0, 0, 1, 1, 0), (0, 0, 1, 0, 1, 1, 0), (1, 0, 1, 0, 1, 0, 1),
 (0, 1, 1, 0, 0, 1, 1), (1, 1, 1, 0, 0, 0, 0), (0, 0, 0, 1, 1, 1, 1),
 (1, 0, 0, 1, 1, 0, 0), (0, 1, 0, 1, 0, 1, 0), (1, 1, 0, 1, 0, 0, 1),
 (0, 0, 1, 1, 0, 0, 1), (1, 0, 1, 1, 0, 1, 0), (0, 1, 1, 1, 1, 0, 0),
 (1, 1, 1, 1, 1, 1, 1)]
```

The minimum distance is perhaps one of the most important properties. Hamming codes always have minimum distance $d = 3$, so they are always single error-correcting.

```
H.minimum_distance()
```

```
3
```

We know that the parity-check matrix and the generator matrix are useful for the construction, description and analysis of linear codes. The Sage method names are just a bit cryptic. Sage has extensive routines for analyzing matrices with elements from different fields, so we perform much of the subsequent analysis of these matrices within Sage.

```
C = H.parity_check_matrix(); C
```

```
[1 0 1 0 1 0 1]
[0 1 1 0 0 1 1]
[0 0 0 1 1 1 1]
```

The generator matrix here in the text has *columns* that are codewords, and linear combinations of the columns (the column space of the matrix) are codewords. In Sage the generator matrix has *rows* that are codewords and the row space of the matrix is the code. So here is another place where we need to mentally translate between a choice made in the text and a choice made by the Sage developers.

```
G = H.generator_matrix(); G
```

```
[1 0 0 0 0 1 1]
[0 1 0 0 1 0 1]
[0 0 1 0 1 1 0]
[0 0 0 1 1 1 1]
```

Here is a partial test that these two matrices are correct, exercising Lemma 8.27. Notice that we need to use the transpose of the generator matrix, for reasons described above.

```
C*G.transpose() == zero_matrix(3, 4)
```

```
True
```

Note that the parity-check may not be canonical and the generator matrix may not be standard. Sage can produce a generator matrix that has a set of columns that forms an identity matrix, though no guarantee is made that these columns are the first columns. (Columns, not rows.) Such a matrix is said to be **systematic**, and the Sage method is `.systematic_generator_matrix()`.

```
H.systematic_generator_matrix()
```

```
[1 0 0 0 0 1 1]
[0 1 0 0 1 0 1]
[0 0 1 0 1 1 0]
[0 0 0 1 1 1 1]
```

Decoding with a Linear Code

We can decode received messages originating from a linear code. Suppose we receive the length 7 binary vector `r`.

```
r = vector(GF(2), [1, 1, 1, 1, 0, 0, 1]); r
```

```
(1, 1, 1, 1, 0, 0, 1)
```

We can recognize that one or more errors has occured, since `r` is not in the code, as the next computation does not yield the zero vector.

```
C*r
```

```
(1, 1, 0)
```

A linear code has a `.decode` method. You may choose from several different algorithms, while the Hamming codes have their own custom algorithm. The default algorithm is syndrome decoding.

```
H.decode_to_code(r)
```

```
(1, 1, 0, 1, 0, 0, 1)
```

So if we are willing to assume that only one error occured (which we might, if the probability of an indivual entry of the vector being in error is very low), then we see that an error occured in the third position.

Remember that it could happen that there was more than just one error. For example, suppose the message was the same as before and errors occurred in the third, fifth and sixth locations.

```
message = vector(GF(2), [1, 1, 0, 1, 0, 0, 1])
errors = vector(GF(2), [0, 0, 1, 0, 1, 1, 0])
received = message + errors
received
```

```
(1, 1, 1, 1, 1, 1, 1)
```

It then appears that we have received a codeword, so we assume no errors at all, and decode incorrectly.

```
H.decode_to_code(received) == message
```

```
False
```

```
H.decode_to_code(received) == received
```

```
True
```

8.8 Sage Exercises

1. Create the (binary) Golay code with the `codes.GolayCode()` constructor. Read the documentation to be sure you build the binary version (not ternary), and do not build the extended version (which is the default).

(a) Use Sage methods to compute the length, dimension and minimum distance of the code.

(b) How many errors can this code detect? How many can it correct?

(c) Find a nonzero codeword and introduce three errors by adding a vector with three 1's (your choice) to create a received message. Show that the message is decoded properly.

(d) Recycle your choices from the previous part, but now add one more error. Does the new received message get decoded properly?

2. One technique for improving the characteristics of a code is to add an overall parity-check bit, much like the lone parity-check bit of the ASCII code described in Example 8.3. Such codes are referred to as the ***extended*** version of the original.

(a) Construct the (binary) Golay code and obtain the parity-check matrix. Use Sage commands to enlarge this matrix to create a new parity check matrix that has an additional overall parity-check bit. You may find the matrix methods `.augment()` and `.stack()` useful, as well as the constructors `zero_vector()` and `ones_matrix()` (remembering that we specify the binary entries as being from the field `GF(2)`.)

Create the extended code by supplying your enlarged parity-check matrix to the `codes.from_parity_check_matrix()` constructor and compute the length, dimension and minimum distance of the extended code.

(b) How are the properties of this new code better? At what cost?

(c) Now create the extended (binary) Golay code with the Sage constructor `codes.GolayCode()` and the correct keyword to obtain the extended version. With luck, the sorted lists of your codewords and Sage's codewords will be equal. If not, the linear code method `.is_permutation_equivalent()` should return `True` to indicate that your code and Sage's are just rearrangements of each other.

3. *Note:* This problem is on holiday (as of Sage 6.7), while some buggy Sage code for the minimum distance of a Hamming code gets sorted out. The `r = 2` case produces an error message and for `r > 5` the computation of the minimum distance has become intolerably slow. So it is a bit harder to make a reasonable conjecture from just 3 cases.

The dual of an (n, k) block code is formed as all the set of all binary vectors which are orthogonal to every vector of the original code. Exercise 8.5.25 describes this construction and asks about some of its properties.

You can construct the dual of a code in Sage with the `.dual_code()` method. Construct the binary Hamming codes, and their duals, with the parameter `r` ranging from 2 to 5, inclusive. Build a table with six columns (perhaps employing the `html.table()` function) that lists r, the length of the codes, the dimensions of the original and the dual, and the minimum distances of the orginal and the dual.

Conjecture formulas for the dimension and minimum distance of the dual of the Hamming code as expressions in the parameter r.

4.   A code with minimum distance d is called ***perfect*** if every possible vector is within Hamming distance $(d-1)/2$ of some codeword. If we expand our notion of geometry to account for the Hamming distance as the metric, then we can speak of a sphere of radius r around a vector (or codeword. For a code of length n, such a sphere will contain

$$1 + \binom{n}{1} + \binom{n}{2} + \cdots + \binom{n}{r}$$

vectors within in it. For a perfect code, the spheres of radius d centered at the codewords of the code will exactly partition the entire set of all possible vectors. (This is the connection that means that coding theory meshes with sphere packing problems.)

A consequence of a code of dimension k being perfect is that

$$2^k \left(\binom{n}{0} + \binom{n}{1} + \binom{n}{2} + \cdots + \binom{n}{\frac{d-1}{2}} \right) = 2^n$$

Conversely, if a code has minimum distance d and the condition above is true, then the code is perfect.

Write a Python function, named `is_perfect()` which accepts a linear code as input and returns `True` or `False`. Demonstrate your function by checking that the (binary) Golay code is perfect, and then use a loop to verify that the (binary) Hamming codes are perfect for all lengths below 32.

<div style="text-align: center;">

$\boxed{9}$

Isomorphisms

</div>

Many groups may appear to be different at first glance, but can be shown to be the same by a simple renaming of the group elements. For example, $\mathbb{Z}_4$ and the subgroup of the circle group $\mathbb{T}$ generated by i can be shown to be the same by demonstrating a one-to-one correspondence between the elements of the two groups and between the group operations. In such a case we say that the groups are isomorphic.

9.1 Definition and Examples

Two groups $(G, \cdot)$ and $(H, \circ)$ are **_isomorphic_** if there exists a one-to-one and onto map $\phi : G \to H$ such that the group operation is preserved; that is,

$$\phi(a \cdot b) = \phi(a) \circ \phi(b)$$

for all a and b in G. If G is isomorphic to H, we write $G \cong H$. The map ϕ is called an **_isomorphism_**.

Example 9.1. To show that $\mathbb{Z}_4 \cong \langle i \rangle$, define a map $\phi : \mathbb{Z}_4 \to \langle i \rangle$ by $\phi(n) = i^n$. We must show that ϕ is bijective and preserves the group operation. The map ϕ is one-to-one and onto because

$$\phi(0) = 1$$
$$\phi(1) = i$$
$$\phi(2) = -1$$
$$\phi(3) = -i.$$

Since

$$\phi(m + n) = i^{m+n} = i^m i^n = \phi(m)\phi(n),$$

the group operation is preserved.

Example 9.2. We can define an isomorphism ϕ from the additive group of real numbers $(\mathbb{R}, +)$ to the multiplicative group of positive real numbers $(\mathbb{R}^+, \cdot)$ with the exponential map; that is,

$$\phi(x + y) = e^{x+y} = e^x e^y = \phi(x)\phi(y).$$

Of course, we must still show that ϕ is one-to-one and onto, but this can be determined using calculus.

Example 9.3. The integers are isomorphic to the subgroup of $\mathbb{Q}^*$ consisting of elements of the form 2^n. Define a map $\phi : \mathbb{Z} \to \mathbb{Q}^*$ by $\phi(n) = 2^n$. Then

$$\phi(m + n) = 2^{m+n} = 2^m 2^n = \phi(m)\phi(n).$$

<div style="text-align: center;">153</div>

By definition the map ϕ is onto the subset $\{2^n : n \in \mathbb{Z}\}$ of $\mathbb{Q}^*$. To show that the map is injective, assume that $m \neq n$. If we can show that $\phi(m) \neq \phi(n)$, then we are done. Suppose that $m > n$ and assume that $\phi(m) = \phi(n)$. Then $2^m = 2^n$ or $2^{m-n} = 1$, which is impossible since $m - n > 0$.

Example 9.4. The groups $\mathbb{Z}_8$ and $\mathbb{Z}_{12}$ cannot be isomorphic since they have different orders; however, it is true that $U(8) \cong U(12)$. We know that

$$U(8) = \{1, 3, 5, 7\}$$
$$U(12) = \{1, 5, 7, 11\}.$$

An isomorphism $\phi : U(8) \to U(12)$ is then given by

$$1 \mapsto 1$$
$$3 \mapsto 5$$
$$5 \mapsto 7$$
$$7 \mapsto 11.$$

The map ϕ is not the only possible isomorphism between these two groups. We could define another isomorphism ψ by $\psi(1) = 1$, $\psi(3) = 11$, $\psi(5) = 5$, $\psi(7) = 7$. In fact, both of these groups are isomorphic to $\mathbb{Z}_2 \times \mathbb{Z}_2$ (see Example 3.28 in Chapter 3).

Example 9.5. Even though S_3 and $\mathbb{Z}_6$ possess the same number of elements, we would suspect that they are not isomorphic, because $\mathbb{Z}_6$ is abelian and S_3 is nonabelian. To demonstrate that this is indeed the case, suppose that $\phi : \mathbb{Z}_6 \to S_3$ is an isomorphism. Let $a, b \in S_3$ be two elements such that $ab \neq ba$. Since ϕ is an isomorphism, there exist elements m and n in $\mathbb{Z}_6$ such that

$$\phi(m) = a \quad \text{and} \quad \phi(n) = b.$$

However,

$$ab = \phi(m)\phi(n) = \phi(m + n) = \phi(n + m) = \phi(n)\phi(m) = ba,$$

which contradicts the fact that a and b do not commute.

Theorem 9.6. *Let $\phi : G \to H$ be an isomorphism of two groups. Then the following statements are true.*

1. *$\phi^{-1} : H \to G$ is an isomorphism.*

2. *$|G| = |H|$.*

3. *If G is abelian, then H is abelian.*

4. *If G is cyclic, then H is cyclic.*

5. *If G has a subgroup of order n, then H has a subgroup of order n.*

PROOF.   Assertions (1) and (2) follow from the fact that ϕ is a bijection. We will prove (3) here and leave the remainder of the theorem to be proved in the exercises.

(3) Suppose that h_1 and h_2 are elements of H. Since ϕ is onto, there exist elements $g_1, g_2 \in G$ such that $\phi(g_1) = h_1$ and $\phi(g_2) = h_2$. Therefore,

$$h_1 h_2 = \phi(g_1)\phi(g_2) = \phi(g_1 g_2) = \phi(g_2 g_1) = \phi(g_2)\phi(g_1) = h_2 h_1. \qquad \square$$

We are now in a position to characterize all cyclic groups.

Theorem 9.7. *All cyclic groups of infinite order are isomorphic to $\mathbb{Z}$.*

PROOF. Let G be a cyclic group with infinite order and suppose that a is a generator of G. Define a map $\phi : \mathbb{Z} \to G$ by $\phi : n \mapsto a^n$. Then

$$\phi(m + n) = a^{m+n} = a^m a^n = \phi(m)\phi(n).$$

To show that ϕ is injective, suppose that m and n are two elements in $\mathbb{Z}$, where $m \neq n$. We can assume that $m > n$. We must show that $a^m \neq a^n$. Let us suppose the contrary; that is, $a^m = a^n$. In this case $a^{m-n} = e$, where $m - n > 0$, which contradicts the fact that a has infinite order. Our map is onto since any element in G can be written as a^n for some integer n and $\phi(n) = a^n$. $\square$

Theorem 9.8. *If G is a cyclic group of order n, then G is isomorphic to $\mathbb{Z}_n$.*

PROOF. Let G be a cyclic group of order n generated by a and define a map $\phi : \mathbb{Z}_n \to G$ by $\phi : k \mapsto a^k$, where $0 \leq k < n$. The proof that ϕ is an isomorphism is one of the end-of-chapter exercises. $\square$

Corollary 9.9. *If G is a group of order p, where p is a prime number, then G is isomorphic to $\mathbb{Z}_p$.*

PROOF. The proof is a direct result of Corollary 6.12. $\square$

The main goal in group theory is to classify all groups; however, it makes sense to consider two groups to be the same if they are isomorphic. We state this result in the following theorem, whose proof is left as an exercise.

Theorem 9.10. *The isomorphism of groups determines an equivalence relation on the class of all groups.*

Hence, we can modify our goal of classifying all groups to classifying all groups **up to isomorphism**; that is, we will consider two groups to be the same if they are isomorphic.

Cayley's Theorem

Cayley proved that if G is a group, it is isomorphic to a group of permutations on some set; hence, every group is a permutation group. Cayley's Theorem is what we call a representation theorem. The aim of representation theory is to find an isomorphism of some group G that we wish to study into a group that we know a great deal about, such as a group of permutations or matrices.

Example 9.11. Consider the group $\mathbb{Z}_3$. The Cayley table for $\mathbb{Z}_3$ is as follows.

+	0	1	2
0	0	1	2
1	1	2	0
2	2	0	1

The addition table of $\mathbb{Z}_3$ suggests that it is the same as the permutation group $G = \{(0), (012), (021)\}$. The isomorphism here is

$$0 \mapsto \begin{pmatrix} 0 & 1 & 2 \\ 0 & 1 & 2 \end{pmatrix} = (0)$$

$$1 \mapsto \begin{pmatrix} 0 & 1 & 2 \\ 1 & 2 & 0 \end{pmatrix} = (012)$$

$$2 \mapsto \begin{pmatrix} 0 & 1 & 2 \\ 2 & 0 & 1 \end{pmatrix} = (021).$$

Theorem 9.12 Cayley. *Every group is isomorphic to a group of permutations.*

PROOF.   Let G be a group. We must find a group of permutations $\overline{G}$ that is isomorphic to G. For any $g \in G$, define a function $\lambda_g : G \to G$ by $\lambda_g(a) = ga$. We claim that λ_g is a permutation of G. To show that λ_g is one-to-one, suppose that $\lambda_g(a) = \lambda_g(b)$. Then

$$ga = \lambda_g(a) = \lambda_g(b) = gb.$$

Hence, $a = b$. To show that λ_g is onto, we must prove that for each $a \in G$, there is a b such that $\lambda_g(b) = a$. Let $b = g^{-1}a$.

Now we are ready to define our group $\overline{G}$. Let

$$\overline{G} = \{\lambda_g : g \in G\}.$$

We must show that $\overline{G}$ is a group under composition of functions and find an isomorphism between G and $\overline{G}$. We have closure under composition of functions since

$$(\lambda_g \circ \lambda_h)(a) = \lambda_g(ha) = gha = \lambda_{gh}(a).$$

Also,

$$\lambda_e(a) = ea = a$$

and

$$(\lambda_{g^{-1}} \circ \lambda_g)(a) = \lambda_{g^{-1}}(ga) = g^{-1}ga = a = \lambda_e(a).$$

We can define an isomorphism from G to $\overline{G}$ by $\phi : g \mapsto \lambda_g$. The group operation is preserved since

$$\phi(gh) = \lambda_{gh} = \lambda_g \lambda_h = \phi(g)\phi(h).$$

It is also one-to-one, because if $\phi(g)(a) = \phi(h)(a)$, then

$$ga = \lambda_g a = \lambda_h a = ha.$$

Hence, $g = h$. That ϕ is onto follows from the fact that $\phi(g) = \lambda_g$ for any $\lambda_g \in \overline{G}$.   □

The isomorphism $g \mapsto \lambda_g$ is known as the ***left regular representation*** of G.

Historical Note

Arthur Cayley was born in England in 1821, though he spent much of the first part of his life in Russia, where his father was a merchant. Cayley was educated at Cambridge, where he took the first Smith's Prize in mathematics. A lawyer for much of his adult life, he wrote several papers in his early twenties before entering the legal profession at the age of 25. While practicing law he continued his mathematical research, writing more than 300 papers during this period of his life. These included some of his best work. In 1863 he left law to become a professor at Cambridge. Cayley wrote more than 900 papers in fields such as group theory, geometry, and linear algebra. His legal knowledge was very valuable to Cambridge; he participated in the writing of many of the university's statutes. Cayley was also one of the people responsible for the admission of women to Cambridge.

9.2 Direct Products

Given two groups G and H, it is possible to construct a new group from the Cartesian product of G and H, $G \times H$. Conversely, given a large group, it is sometimes possible to decompose the group; that is, a group is sometimes isomorphic to the direct product of two smaller groups. Rather than studying a large group G, it is often easier to study the component groups of G.

External Direct Products

If $(G, \cdot)$ and $(H, \circ)$ are groups, then we can make the Cartesian product of G and H into a new group. As a set, our group is just the ordered pairs $(g, h) \in G \times H$ where $g \in G$ and $h \in H$. We can define a binary operation on $G \times H$ by

$$(g_1, h_1)(g_2, h_2) = (g_1 \cdot g_2, h_1 \circ h_2);$$

that is, we just multiply elements in the first coordinate as we do in G and elements in the second coordinate as we do in H. We have specified the particular operations $\cdot$ and $\circ$ in each group here for the sake of clarity; we usually just write $(g_1, h_1)(g_2, h_2) = (g_1 g_2, h_1 h_2)$.

Proposition 9.13. *Let G and H be groups. The set $G \times H$ is a group under the operation* $(g_1, h_1)(g_2, h_2) = (g_1 g_2, h_1 h_2)$ *where* $g_1, g_2 \in G$ *and* $h_1, h_2 \in H$.

PROOF. Clearly the binary operation defined above is closed. If e_G and e_H are the identities of the groups G and H respectively, then (e_G, e_H) is the identity of $G \times H$. The inverse of $(g, h) \in G \times H$ is (g^{-1}, h^{-1}). The fact that the operation is associative follows directly from the associativity of G and H. $\square$

Example 9.14. Let $\mathbb{R}$ be the group of real numbers under addition. The Cartesian product of $\mathbb{R}$ with itself, $\mathbb{R} \times \mathbb{R} = \mathbb{R}^2$, is also a group, in which the group operation is just addition in each coordinate; that is, $(a, b) + (c, d) = (a + c, b + d)$. The identity is $(0, 0)$ and the inverse of (a, b) is $(-a, -b)$.

Example 9.15. Consider

$$\mathbb{Z}_2 \times \mathbb{Z}_2 = \{(0, 0), (0, 1), (1, 0), (1, 1)\}.$$

Although $\mathbb{Z}_2 \times \mathbb{Z}_2$ and $\mathbb{Z}_4$ both contain four elements, they are not isomorphic. Every element (a, b) in $\mathbb{Z}_2 \times \mathbb{Z}_2$ has order 2, since $(a, b) + (a, b) = (0, 0)$; however, $\mathbb{Z}_4$ is cyclic.

The group $G \times H$ is called the ***external direct product*** of G and H. Notice that there is nothing special about the fact that we have used only two groups to build a new group. The direct product

$$\prod_{i=1}^{n} G_i = G_1 \times G_2 \times \cdots \times G_n$$

of the groups $G_1, G_2, \ldots, G_n$ is defined in exactly the same manner. If $G = G_1 = G_2 = \cdots = G_n$, we often write G^n instead of $G_1 \times G_2 \times \cdots \times G_n$.

Example 9.16. The group $\mathbb{Z}_2^n$, considered as a set, is just the set of all binary n-tuples. The group operation is the "exclusive or" of two binary n-tuples. For example,

$$(01011101) + (01001011) = (00010110).$$

This group is important in coding theory, in cryptography, and in many areas of computer science.

Theorem 9.17. *Let $(g, h) \in G \times H$. If g and h have finite orders r and s respectively, then the order of (g, h) in $G \times H$ is the least common multiple of r and s.*

PROOF.   Suppose that m is the least common multiple of r and s and let $n = |(g, h)|$. Then

$$(g, h)^m = (g^m, h^m) = (e_G, e_H)$$
$$(g^n, h^n) = (g, h)^n = (e_G, e_H).$$

Hence, n must divide m, and $n \leq m$. However, by the second equation, both r and s must divide n; therefore, n is a common multiple of r and s. Since m is the *least common multiple* of r and s, $m \leq n$. Consequently, m must be equal to n. $\square$

Corollary 9.18. *Let $(g_1, \ldots, g_n) \in \prod G_i$. If g_i has finite order r_i in G_i, then the order of $(g_1, \ldots, g_n)$ in $\prod G_i$ is the least common multiple of $r_1, \ldots, r_n$.*

Example 9.19. Let $(8, 56) \in \mathbb{Z}_{12} \times \mathbb{Z}_{60}$. Since $\gcd(8, 12) = 4$, the order of 8 is $12/4 = 3$ in $\mathbb{Z}_{12}$. Similarly, the order of 56 in $\mathbb{Z}_{60}$ is 15. The least common multiple of 3 and 15 is 15; hence, $(8, 56)$ has order 15 in $\mathbb{Z}_{12} \times \mathbb{Z}_{60}$.

Example 9.20. The group $\mathbb{Z}_2 \times \mathbb{Z}_3$ consists of the pairs

$$(0, 0), \qquad (0, 1), \qquad (0, 2), \qquad (1, 0), \qquad (1, 1), \qquad (1, 2).$$

In this case, unlike that of $\mathbb{Z}_2 \times \mathbb{Z}_2$ and $\mathbb{Z}_4$, it is true that $\mathbb{Z}_2 \times \mathbb{Z}_3 \cong \mathbb{Z}_6$. We need only show that $\mathbb{Z}_2 \times \mathbb{Z}_3$ is cyclic. It is easy to see that $(1, 1)$ is a generator for $\mathbb{Z}_2 \times \mathbb{Z}_3$.

The next theorem tells us exactly when the direct product of two cyclic groups is cyclic.

Theorem 9.21. *The group $\mathbb{Z}_m \times \mathbb{Z}_n$ is isomorphic to $\mathbb{Z}_{mn}$ if and only if $\gcd(m, n) = 1$.*

PROOF.   We will first show that if $\mathbb{Z}_m \times \mathbb{Z}_n \cong \mathbb{Z}_{mn}$, then $\gcd(m, n) = 1$. We will prove the contrapositive; that is, we will show that if $\gcd(m, n) = d > 1$, then $\mathbb{Z}_m \times \mathbb{Z}_n$ cannot be cyclic. Notice that mn/d is divisible by both m and n; hence, for any element $(a, b) \in \mathbb{Z}_m \times \mathbb{Z}_n$,

$$\underbrace{(a, b) + (a, b) + \cdots + (a, b)}_{mn/d \text{ times}} = (0, 0).$$

Therefore, no (a, b) can generate all of $\mathbb{Z}_m \times \mathbb{Z}_n$.

The converse follows directly from Theorem 9.17 since $\operatorname{lcm}(m, n) = mn$ if and only if $\gcd(m, n) = 1$. $\square$

Corollary 9.22. *Let $n_1, \ldots, n_k$ be positive integers. Then*

$$\prod_{i=1}^{k} \mathbb{Z}_{n_i} \cong \mathbb{Z}_{n_1 \cdots n_k}$$

if and only if $\gcd(n_i, n_j) = 1$ for $i \neq j$.

Corollary 9.23. *If*

$$m = p_1^{e_1} \cdots p_k^{e_k},$$

where the p_is are distinct primes, then

$$\mathbb{Z}_m \cong \mathbb{Z}_{p_1^{e_1}} \times \cdots \times \mathbb{Z}_{p_k^{e_k}}.$$

PROOF. Since the greatest common divisor of $p_i^{e_i}$ and $p_j^{e_j}$ is 1 for $i \neq j$, the proof follows from Corollary 9.22. □

In Chapter 13, we will prove that all finite abelian groups are isomorphic to direct products of the form

$$\mathbb{Z}_{p_1^{e_1}} \times \cdots \times \mathbb{Z}_{p_k^{e_k}}$$

where $p_1, \ldots, p_k$ are (not necessarily distinct) primes.

Internal Direct Products

The external direct product of two groups builds a large group out of two smaller groups. We would like to be able to reverse this process and conveniently break down a group into its direct product components; that is, we would like to be able to say when a group is isomorphic to the direct product of two of its subgroups.

Let G be a group with subgroups H and K satisfying the following conditions.

- $G = HK = \{hk : h \in H, k \in K\}$;

- $H \cap K = \{e\}$;

- $hk = kh$ for all $k \in K$ and $h \in H$.

Then G is the ***internal direct product*** of H and K.

Example 9.24. The group $U(8)$ is the internal direct product of

$$H = \{1, 3\} \quad \text{and} \quad K = \{1, 5\}.$$

Example 9.25. The dihedral group D_6 is an internal direct product of its two subgroups

$$H = \{\text{id}, r^3\} \quad \text{and} \quad K = \{\text{id}, r^2, r^4, s, r^2 s, r^4 s\}.$$

It can easily be shown that $K \cong S_3$; consequently, $D_6 \cong \mathbb{Z}_2 \times S_3$.

Example 9.26. Not every group can be written as the internal direct product of two of its proper subgroups. If the group S_3 were an internal direct product of its proper subgroups H and K, then one of the subgroups, say H, would have to have order 3. In this case H is the subgroup $\{(1), (123), (132)\}$. The subgroup K must have order 2, but no matter which subgroup we choose for K, the condition that $hk = kh$ will never be satisfied for $h \in H$ and $k \in K$.

Theorem 9.27. *Let G be the internal direct product of subgroups H and K. Then G is isomorphic to $H \times K$.*

PROOF. Since G is an internal direct product, we can write any element $g \in G$ as $g = hk$ for some $h \in H$ and some $k \in K$. Define a map $\phi : G \to H \times K$ by $\phi(g) = (h, k)$.

The first problem that we must face is to show that ϕ is a well-defined map; that is, we must show that h and k are uniquely determined by g. Suppose that $g = hk = h'k'$. Then $h^{-1}h' = k(k')^{-1}$ is in both H and K, so it must be the identity. Therefore, $h = h'$ and $k = k'$, which proves that ϕ is, indeed, well-defined.

To show that ϕ preserves the group operation, let $g_1 = h_1 k_1$ and $g_2 = h_2 k_2$ and observe that

$$\phi(g_1 g_2) = \phi(h_1 k_1 h_2 k_2)$$
$$= \phi(h_1 h_2 k_1 k_2)$$

$$= (h_1 h_2, k_1 k_2)$$
$$= (h_1, k_1)(h_2, k_2)$$
$$= \phi(g_1)\phi(g_2).$$

We will leave the proof that ϕ is one-to-one and onto as an exercise. $\qquad\qquad$ $\square$

Example 9.28. The group $\mathbb{Z}_6$ is an internal direct product isomorphic to $\{0, 2, 4\} \times \{0, 3\}$.

We can extend the definition of an internal direct product of G to a collection of subgroups $H_1, H_2, \ldots, H_n$ of G, by requiring that

- $G = H_1 H_2 \cdots H_n = \{h_1 h_2 \cdots h_n : h_i \in H_i\}$;

- $H_i \cap \langle \cup_{j \neq i} H_j \rangle = \{e\}$;

- $h_i h_j = h_j h_i$ for all $h_i \in H_i$ and $h_j \in H_j$.

We will leave the proof of the following theorem as an exercise.

Theorem 9.29. *Let G be the internal direct product of subgroups H_i, where $i = 1, 2, \ldots, n$. Then G is isomorphic to $\prod_i H_i$.*

9.3   Exercises

1. Prove that $\mathbb{Z} \cong n\mathbb{Z}$ for $n \neq 0$.

2. Prove that $\mathbb{C}^*$ is isomorphic to the subgroup of $GL_2(\mathbb{R})$ consisting of matrices of the form

$$\begin{pmatrix} a & b \\ -b & a \end{pmatrix}.$$

3. Prove or disprove: $U(8) \cong \mathbb{Z}_4$.

4. Prove that $U(8)$ is isomorphic to the group of matrices

$$\begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}, \begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix}, \begin{pmatrix} -1 & 0 \\ 0 & 1 \end{pmatrix}, \begin{pmatrix} -1 & 0 \\ 0 & -1 \end{pmatrix}.$$

5. Show that $U(5)$ is isomorphic to $U(10)$, but $U(12)$ is not.

6. Show that the nth roots of unity are isomorphic to $\mathbb{Z}_n$.

7. Show that any cyclic group of order n is isomorphic to $\mathbb{Z}_n$.

8. Prove that $\mathbb{Q}$ is not isomorphic to $\mathbb{Z}$.

9. Let $G = \mathbb{R} \setminus \{-1\}$ and define a binary operation on G by

$$a * b = a + b + ab.$$

Prove that G is a group under this operation. Show that $(G, *)$ is isomorphic to the multiplicative group of nonzero real numbers.

10. Show that the matrices

$$\begin{pmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1 \end{pmatrix} \quad \begin{pmatrix} 1 & 0 & 0 \\ 0 & 0 & 1 \\ 0 & 1 & 0 \end{pmatrix} \quad \begin{pmatrix} 0 & 1 & 0 \\ 1 & 0 & 0 \\ 0 & 0 & 1 \end{pmatrix}$$

$$\begin{pmatrix} 0 & 0 & 1 \\ 1 & 0 & 0 \\ 0 & 1 & 0 \end{pmatrix} \quad \begin{pmatrix} 0 & 0 & 1 \\ 0 & 1 & 0 \\ 1 & 0 & 0 \end{pmatrix} \quad \begin{pmatrix} 0 & 1 & 0 \\ 0 & 0 & 1 \\ 1 & 0 & 0 \end{pmatrix}$$

form a group. Find an isomorphism of G with a more familiar group of order 6.

11.   Find five non-isomorphic groups of order 8.

12.   Prove S_4 is not isomorphic to D_{12}.

13.   Let $\omega = \operatorname{cis}(2\pi/n)$ be a primitive nth root of unity. Prove that the matrices

$$A = \begin{pmatrix} \omega & 0 \\ 0 & \omega^{-1} \end{pmatrix} \quad \text{and} \quad B = \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix}$$

generate a multiplicative group isomorphic to D_n.

14.   Show that the set of all matrices of the form

$$\begin{pmatrix} \pm 1 & k \\ 0 & 1 \end{pmatrix},$$

is a group isomorphic to D_n, where all entries in the matrix are in $\mathbb{Z}_n$.

15.   List all of the elements of $\mathbb{Z}_4 \times \mathbb{Z}_2$.

16.   Find the order of each of the following elements.

(a) $(3, 4)$ in $\mathbb{Z}_4 \times \mathbb{Z}_6$

(b) $(6, 15, 4)$ in $\mathbb{Z}_{30} \times \mathbb{Z}_{45} \times \mathbb{Z}_{24}$

(c) $(5, 10, 15)$ in $\mathbb{Z}_{25} \times \mathbb{Z}_{25} \times \mathbb{Z}_{25}$

(d) $(8, 8, 8)$ in $\mathbb{Z}_{10} \times \mathbb{Z}_{24} \times \mathbb{Z}_{80}$

17.   Prove that D_4 cannot be the internal direct product of two of its proper subgroups.

18.   Prove that the subgroup of $\mathbb{Q}^*$ consisting of elements of the form $2^m 3^n$ for $m, n \in \mathbb{Z}$ is an internal direct product isomorphic to $\mathbb{Z} \times \mathbb{Z}$.

19.   Prove that $S_3 \times \mathbb{Z}_2$ is isomorphic to D_6. Can you make a conjecture about D_{2n}? Prove your conjecture.

20.   Prove or disprove: Every abelian group of order divisible by 3 contains a subgroup of order 3.

21.   Prove or disprove: Every nonabelian group of order divisible by 6 contains a subgroup of order 6.

22.   Let G be a group of order 20. If G has subgroups H and K of orders 4 and 5 respectively such that $hk = kh$ for all $h \in H$ and $k \in K$, prove that G is the internal direct product of H and K.

23.   Prove or disprove the following assertion. Let G, H, and K be groups. If $G \times K \cong H \times K$, then $G \cong H$.

24.   Prove or disprove: There is a noncyclic abelian group of order 51.

25.   Prove or disprove: There is a noncyclic abelian group of order 52.

26.   Let $\phi : G \to H$ be a group isomorphism. Show that $\phi(x) = e_H$ if and only if $x = e_G$, where e_G and e_H are the identities of G and H, respectively.

27.   Let $G \cong H$. Show that if G is cyclic, then so is H.

28.   Prove that any group G of order p, p prime, must be isomorphic to $\mathbb{Z}_p$.

29.   Show that S_n is isomorphic to a subgroup of A_{n+2}.

30.   Prove that D_n is isomorphic to a subgroup of S_n.

31.   Let $\phi : G_1 \to G_2$ and $\psi : G_2 \to G_3$ be isomorphisms. Show that ϕ^{-1} and $\psi \circ \phi$ are both isomorphisms. Using these results, show that the isomorphism of groups determines an equivalence relation on the class of all groups.

32.   Prove $U(5) \cong \mathbb{Z}_4$. Can you generalize this result for $U(p)$, where p is prime?

33.   Write out the permutations associated with each element of S_3 in the proof of Cayley's Theorem.

34.   An ***automorphism*** of a group G is an isomorphism with itself. Prove that complex conjugation is an automorphism of the additive group of complex numbers; that is, show that the map $\phi(a + bi) = a - bi$ is an isomorphism from $\mathbb{C}$ to $\mathbb{C}$.

35.   Prove that $a + ib \mapsto a - ib$ is an automorphism of $\mathbb{C}^*$.

36.   Prove that $A \mapsto B^{-1}AB$ is an automorphism of $SL_2(\mathbb{R})$ for all B in $GL_2(\mathbb{R})$.

37.   We will denote the set of all automorphisms of G by $\mathrm{Aut}(G)$. Prove that $\mathrm{Aut}(G)$ is a subgroup of S_G, the group of permutations of G.

38.   Find $\mathrm{Aut}(\mathbb{Z}_6)$.

39.   Find $\mathrm{Aut}(\mathbb{Z})$.

40.   Find two nonisomorphic groups G and H such that $\mathrm{Aut}(G) \cong \mathrm{Aut}(H)$.

41.   Let G be a group and $g \in G$. Define a map $i_g : G \to G$ by $i_g(x) = gxg^{-1}$. Prove that i_g defines an automorphism of G. Such an automorphism is called an ***inner automorphism***. The set of all inner automorphisms is denoted by $\mathrm{Inn}(G)$.

42.   Prove that $\mathrm{Inn}(G)$ is a subgroup of $\mathrm{Aut}(G)$.

43.   What are the inner automorphisms of the quaternion group Q_8? Is $\mathrm{Inn}(G) = \mathrm{Aut}(G)$ in this case?

44.   Let G be a group and $g \in G$. Define maps $\lambda_g : G \to G$ and $\rho_g : G \to G$ by $\lambda_g(x) = gx$ and $\rho_g(x) = xg^{-1}$. Show that $i_g = \rho_g \circ \lambda_g$ is an automorphism of G. The isomorphism $g \mapsto \rho_g$ is called the ***right regular representation*** of G.

45.   Let G be the internal direct product of subgroups H and K. Show that the map $\phi : G \to H \times K$ defined by $\phi(g) = (h, k)$ for $g = hk$, where $h \in H$ and $k \in K$, is one-to-one and onto.

46.   Let G and H be isomorphic groups. If G has a subgroup of order n, prove that H must also have a subgroup of order n.

47.   If $G \cong \overline{G}$ and $H \cong \overline{H}$, show that $G \times H \cong \overline{G} \times \overline{H}$.

48.   Prove that $G \times H$ is isomorphic to $H \times G$.

49.   Let $n_1, \ldots, n_k$ be positive integers. Show that

$$\prod_{i=1}^{k} \mathbb{Z}_{n_i} \cong \mathbb{Z}_{n_1 \cdots n_k}$$

if and only if $\gcd(n_i, n_j) = 1$ for $i \neq j$.

50.   Prove that $A \times B$ is abelian if and only if A and B are abelian.

51.   If G is the internal direct product of $H_1, H_2, \ldots, H_n$, prove that G is isomorphic to $\prod_i H_i$.

52.   Let H_1 and H_2 be subgroups of G_1 and G_2, respectively. Prove that $H_1 \times H_2$ is a subgroup of $G_1 \times G_2$.

53. Let $m, n \in \mathbb{Z}$. Prove that $\langle m, n \rangle = \langle d \rangle$ if and only if $d = \gcd(m, n)$.

54. Let $m, n \in \mathbb{Z}$. Prove that $\langle m \rangle \cap \langle n \rangle = \langle l \rangle$ if and only if $l = \operatorname{lcm}(m, n)$.

55. **Groups of order** $2p$. In this series of exercises we will classify all groups of order $2p$, where p is an odd prime.

(a) Assume G is a group of order $2p$, where p is an odd prime. If $a \in G$, show that a must have order 1, 2, p, or $2p$.

(b) Suppose that G has an element of order $2p$. Prove that G is isomorphic to $\mathbb{Z}_{2p}$. Hence, G is cyclic.

(c) Suppose that G does not contain an element of order $2p$. Show that G must contain an element of order p. *Hint*: Assume that G does not contain an element of order p.

(d) Suppose that G does not contain an element of order $2p$. Show that G must contain an element of order 2.

(e) Let P be a subgroup of G with order p and $y \in G$ have order 2. Show that $yP = Py$.

(f) Suppose that G does not contain an element of order $2p$ and $P = \langle z \rangle$ is a subgroup of order p generated by z. If y is an element of order 2, then $yz = z^k y$ for some $2 \le k < p$.

(g) Suppose that G does not contain an element of order $2p$. Prove that G is not abelian.

(h) Suppose that G does not contain an element of order $2p$ and $P = \langle z \rangle$ is a subgroup of order p generated by z and y is an element of order 2. Show that we can list the elements of G as $\{z^i y^j \mid 0 \le i < p, 0 \le j < 2\}$.

(i) Suppose that G does not contain an element of order $2p$ and $P = \langle z \rangle$ is a subgroup of order p generated by z and y is an element of order 2. Prove that the product $(z^i y^j)(z^r y^s)$ can be expressed as a uniquely as $z^m y^n$ for some non negative integers m, n. Thus, conclude that there is only one possibility for a non-abelian group of order $2p$, it must therefore be the one we have seen already, the dihedral group.

9.4 Sage

Sage has limited support for actually creating isomorphisms, though it is possible. However, there is excellent support for determining if two permutation groups are isomorphic. This will allow us to begin a little project to locate *all* of the groups of order less than 16 in Sage's permutation groups.

Isomorphism Testing

If `G` and `H` are two permutation groups, then the command `G.is_isomorphic(H)` will return `True` or `False` as the two groups are, or are not, isomorphic. Since "isomorpic to" is an equivalence relation by Theorem 9.10, it does not matter which group plays the role of `G` and which plays the role of `H`.

So we have a few more examples to work with, let us introduce the Sage command that creates an external direct product. If `G` and `H` are two permutation groups, then the command `direct_product_permgroups([G,H])` will return the external direct product as a

new permutation group. Notice that this is a function (not a method) and the input is a list. Rather than just combining two groups in the list, any number of groups can be supplied. We illustrate isomorphism testing and direct products in the context of Theorem 9.21, which is an equivalence, so tells us *exactly* when we have isomorphic groups. We use cyclic permutation groups as stand-ins for $\mathbb{Z}_n$ by Theorem 9.8.

First, two isomorphic groups.

```
m = 12
n = 7
gcd(m, n)
```

1

```
G = CyclicPermutationGroup(m)
H = CyclicPermutationGroup(n)
dp = direct_product_permgroups([G, H])
K = CyclicPermutationGroup(m*n)
K.is_isomorphic(dp)
```

True

Now, two non-isomorphic groups.

```
m = 15
n = 21
gcd(m, n)
```

3

```
G = CyclicPermutationGroup(m)
H = CyclicPermutationGroup(n)
dp = direct_product_permgroups([G, H])
K = CyclicPermutationGroup(m*n)
K.is_isomorphic(dp)
```

False

Notice how the simple computation of a greatest common divisor predicts the incredibly complicated computation of determining if two groups are isomorphic. This is a nice illustration of the power of mathematics, replacing a difficult problem (group isomorphism) by a simple one (factoring and divisibility of integers). Let us build one more direct product of cyclic groups, but with three groups, each with orders that are pairwise relatively prime.

If you try the following with larger parameters you may get an error (`database_gap`).

```
m = 6
n = 5
r = 7
G = CyclicPermutationGroup(m)
H = CyclicPermutationGroup(n)
L = CyclicPermutationGroup(r)
dp = direct_product_permgroups([G, H, L])
K = CyclicPermutationGroup(m*n*r)
K.is_isomorphic(dp)
```

True

Classifying Finite Groups

Once we understand isomorphic groups as being the "same", or "fundamentally no different," or "structurally identical," then it is natural to ask how many "really different" finite groups there are. Corollary 9.9 gives a partial answer: for each prime there is just one finite group, with $\mathbb{Z}_p$ as a concrete manifestation.

Let us embark on a quest to find all the groups of order less than 16 in Sage as permutation groups. For prime orders $1, 2, 3, 5, 7, 11$ and 13 we know there is really just one group each, and we can realize them all:

```
[CyclicPermutationGroup(p) for p in [1, 2, 3, 5, 7, 11, 13]]
```

```
[Cyclic group of order 1 as a permutation group,
 Cyclic group of order 2 as a permutation group,
 Cyclic group of order 3 as a permutation group,
 Cyclic group of order 5 as a permutation group,
 Cyclic group of order 7 as a permutation group,
 Cyclic group of order 11 as a permutation group,
 Cyclic group of order 13 as a permutation group]
```

So now our smallest unknown case is order 4. Sage knows at least three such groups, and we can use Sage to check if any pair is isomorphic. Notice that since "isomorphic to" is an equivalence relation, and hence a transitive relation, the two tests below are sufficient.

```
G = CyclicPermutationGroup(4)
H = KleinFourGroup()
T1 = CyclicPermutationGroup(2)
T2 = CyclicPermutationGroup(2)
K = direct_product_permgroups([T1, T2])
G.is_isomorphic(H)
```

```
False
```

```
H.is_isomorphic(K)
```

```
True
```

So we have at least two different groups: $\mathbb{Z}_4$ and $\mathbb{Z}_2 \times \mathbb{Z}_2$, with the latter also known as the Klein 4-group. Sage will not be able to tell us if we have a *complete* list — this will always require theoretical results like Theorem 9.10. We will shortly have a more general result that handles the case of order 4, but right now, a careful analysis (by hand) of the possibilities for the Cayley table of a group of order 4 should lead you to the two possibilities above as the only possibilities. Try to deduce what the Cayley table of an order 4 group should look like, since you know about identity elements, inverses and cancellation.

We have seen at least two groups of order 6 (next on our list of non-prime orders). One is abelian and one is not, so we do not need Sage to tell us they are structurally different. But let us do it anyway.

```
G = CyclicPermutationGroup(6)
H = SymmetricGroup(3)
G.is_isomorphic(H)
```

```
False
```

Is that all? There is $\mathbb{Z}_3 \times \mathbb{Z}_2$, but that is just $\mathbb{Z}_6$ since 2 and 3 are relatively prime. The dihedral group, D_3, all symmetries of a triangle, is just S_3, the symmetric group on 3 symbols.

```
G = DihedralGroup(3)
H = SymmetricGroup(3)
G.is_isomorphic(H)
```

```
True
```

Exercise 9.3.55 from this section classifies all groups of order $2p$, where p is a prime. Such a group is either cyclic or a dihedral group. So the two groups above, $\mathbb{Z}_6$ and D_3, are the complete list of groups of order 6.

By this general result, in addition to order 6, we also know the complete lists of groups of orders 10 and 14. To Be Continued.

Internal Direct Products

An internal direct product is a statement about subgroups of a single group, together with a theorem that links them to an external direct product. We will work an example here that will illustrate the nature of an internal direct product.

Given an integer n, the set of positive integers less than n, and relatively prime to n forms a group under multiplication mod n. We will work in the set Integers(n) where we can add *and* multiply, but we want to stay strictly with multiplication only.

First we build the subgroup itself. Notice how we must convert x into an integer (an element of ZZ) so that the greatest common divisor computation performs correctly.

```
Z36 = Integers(36)
U = [x for x in Z36 if gcd(ZZ(x), 36) == 1]
U
```

```
[1, 5, 7, 11, 13, 17, 19, 23, 25, 29, 31, 35]
```

So we have a group of order 12. We are going to try to find a subgroup of order 6 and a subgroup of order 2 to form the internal direct product, and we will restrict our search initially to cyclic subgroups of order 6. Sage has a method that will give the order of each of these elements, relative to multiplication, so let us examine those next.

```
[x.multiplicative_order() for x in U]
```

```
[1, 6, 6, 6, 3, 2, 2, 6, 3, 6, 6, 2]
```

We have many choices for generators of a cyclic subgroup of order 6 and for a cyclic subgroup of order 2. Of course, some of the choices for a generator of the subgroup of order 6 will generate the same subgroup. Can you tell, just by counting, how many subgroups of order 6 there are? We are going to pick the first element of order 6, and the last element of order 2, for no particular reason. After your work through this once, we encourage you to try other choices to understand why some choices lead to an internal direct product and some do not. Notice that we choose the elements from the list U so that they are sure to be elements of Z36 and behave properly when multiplied.

```
a = U[1]
A = [a^i for i in srange(6)]
A
```

```
[1, 5, 25, 17, 13, 29]
```

```
b = U[11]
B = [b^i for i in srange(2)]
B
```

```
[1, 35]
```

So `A` and `B` are two cyclic subgroups. Notice that their intersection is the identity element, one of our requirements for an internal direct product. So this is a good start.

```
[x for x in A if x in B]
```

```
[1]
```

`Z36` is an abelian group, thus the condition on all products commuting will hold, but we illustrate the Sage commands that will check this in a non-abelian situation.

```
all([x*y == y*x for x in A for y in B])
```

```
True
```

Finally, we need to check that by forming products with elements from `A` and `B` we create the entire group. Sorting the resulting list will make a check easier for us visually, and is required if we want Sage to do the check.

```
T = sorted([x*y for x in A for y in B])
T
```

```
[1, 5, 7, 11, 13, 17, 19, 23, 25, 29, 31, 35]
```

```
T == U
```

```
True
```

That's it. We now condense all this information into the statement that "`U` is the internal direct product of `A` and `B`." By Theorem 9.27, we see that `U` is isomorphic to a product of a cyclic group of order 6 and a cyclic group of order 2. So in a very real sense, `U` is no more or less complicated than $\mathbb{Z}_6 \times \mathbb{Z}_2$, which is in turn isomorphic to $\mathbb{Z}_3 \times \mathbb{Z}_2 \times \mathbb{Z}_2$. So we totally understand the "structure" of `U`. For example, we can see that `U` is not cyclic, since when written as a product of cyclic groups, the two orders are not relatively prime. The final expression of `U` suggests you could find three cyclic subgroups of `U`, with orders 3, 2 and 2, so that `U` is an internal direct product of the three subgroups.

9.5 Sage Exercises

1. This exercise is about putting Cayley's Theorem into practice. First, read and study the theorem. Realize that this result by itself is primarily of theoretical interest, but with some more theory we could get into some subtler aspects of this (a subject known as "representation theory").

You should create these representations mostly with pencil-and-paper work, using Sage as a fancy calculator and assistant. You do not need to include all these computations in your worksheet. Build the requested group representations and then include enough verifications in Sage to prove that that your representation correctly represents the group.

Begin by building a permutation representation of the quaternions, Q. There are eight elements in Q ($\pm 1, \pm I, \pm J, \pm K$), so you will be constructing a subgroup of S_8. For each $g \in Q$ form the function T_g, defined as $T_g(x) = xg$. Notice that this definition is the "reverse" of that given in the text. This is because Sage composes permutations left-to-right, while your text composes right-to-left. To create the permutations T_g, the two-line version of writing permutations could be very useful as an intermediate step. You will probably want to "code" each element of Q with an integer in $\{1, 2, \ldots, 8\}$.

One such representation is included in Sage as `QuaternionGroup()` — your answer should look very similar, but perhaps not identical. Do not submit your answer for a representation

of the quaternions, but I strongly suggest working this particular group representation until you are sure you have it right — the problems below might be very difficult otherwise. You can use Sage's `.is_isomorphic()` method to check if your representations are correct. However, do not use this as a substitute for the part of each question that asks you to investigate properties of your representation towards this end.

(a) Build the permutation representation of $\mathbb{Z}_2 \times \mathbb{Z}_4$ described in Cayley's Theorem. (Remember that this group is additive, while the theorem uses multiplicative notation.) Include the representation of *each* of the 8 elements in your submitted work. Then construct the permutation group as a subgroup of a full symmetric group that is generated by exactly two of the eight elements you have already constructed. Hint: which two elements of $\mathbb{Z}_2 \times \mathbb{Z}_4$ might you use to generate all of $\mathbb{Z}_2 \times \mathbb{Z}_4$? Use commands in Sage to investigate various properties of your permutation group, other than just `.list()`, to provide evidence that your subgroup is correct — include these in your submitted worksheet.

(b) Build a permutation representation of $U(24)$, the group of units mod 24. Again, list a representation of *each* element in your submitted work. Then construct the group as a subgroup of a full symmetric group created with three generators. To determine these three generators, you will likely need to understand $U(24)$ as an internal direct product. Use commands in Sage to investigate various properties of your group, other than just `.list()`, to provide evidence that your subgroup is correct — include these in your submitted worksheet.

2.   Consider the symmetries of a 10-gon, D_{10} in your text, `DihedralGroup(10)` in Sage. Presume that the vertices of the 10-gon have been labeled 1 through 10 in order. Identify the permutation that is a 180 degree rotation and use it to generate a subgroup R of order 2. Then identify the permutation that is a 72 degree rotation, and any one of the ten permutations that are a reflection of the 10-gon about a line. Use these latter two permutations to generate a subgroup S of order 10. Use Sage to verify that the full dihedral group is the internal direct product of the subgroups R and S by checking the conditions in the definition of an internal direct product.

We have a theorem which says that if a group is an internal direct product, then it is isomorphic to some external direct product. Understand that this does not mean that you can use the converse in this problem. In other words, establishing an isomorphism of G with an external direct product *does not prove* that G is an internal direct product.

<div style="text-align: center">

10

</div>

Normal Subgroups and Factor Groups

If H is a subgroup of a group G, then right cosets are not always the same as left cosets; that is, it is not always the case that $gH = Hg$ for all $g \in G$. The subgroups for which this property holds play a critical role in group theory—they allow for the construction of a new class of groups, called factor or quotient groups. Factor groups may be studied directly or by using homomorphisms, a generalization of isomorphisms. We will study homomorphisms in Chapter 11.

10.1 Factor Groups and Normal Subgroups

Normal Subgroups

A subgroup H of a group G is **normal** in G if $gH = Hg$ for all $g \in G$. That is, a normal subgroup of a group G is one in which the right and left cosets are precisely the same.

Example 10.1. Let G be an abelian group. Every subgroup H of G is a normal subgroup. Since $gh = hg$ for all $g \in G$ and $h \in H$, it will always be the case that $gH = Hg$.

Example 10.2. Let H be the subgroup of S_3 consisting of elements (1) and (12). Since

$$(123)H = \{(123), (13)\} \quad \text{and} \quad H(123) = \{(123), (23)\},$$

H cannot be a normal subgroup of S_3. However, the subgroup N, consisting of the permutations (1), (123), and (132), is normal since the cosets of N are

$$N = \{(1), (123), (132)\}$$
$$(12)N = N(12) = \{(12), (13), (23)\}.$$

The following theorem is fundamental to our understanding of normal subgroups.

Theorem 10.3. *Let G be a group and N be a subgroup of G. Then the following statements are equivalent.*

1. *The subgroup N is normal in G.*

2. *For all $g \in G$, $gNg^{-1} \subset N$.*

3. *For all $g \in G$, $gNg^{-1} = N$.*

PROOF.   (1) $\Rightarrow$ (2). Since N is normal in G, $gN = Ng$ for all $g \in G$. Hence, for a given $g \in G$ and $n \in N$, there exists an n' in N such that $gn = n'g$. Therefore, $gng^{-1} = n' \in N$ or $gNg^{-1} \subset N$.

(2) $\Rightarrow$ (3). Let $g \in G$. Since $gNg^{-1} \subset N$, we need only show $N \subset gNg^{-1}$. For $n \in N$, $g^{-1}ng = g^{-1}n(g^{-1})^{-1} \in N$. Hence, $g^{-1}ng = n'$ for some $n' \in N$. Therefore, $n = gn'g^{-1}$ is in gNg^{-1}.

(3) $\Rightarrow$ (1). Suppose that $gNg^{-1} = N$ for all $g \in G$. Then for any $n \in N$ there exists an $n' \in N$ such that $gng^{-1} = n'$. Consequently, $gn = n'g$ or $gN \subset Ng$. Similarly, $Ng \subset gN$. $\qquad\square$

Factor Groups

If N is a normal subgroup of a group G, then the cosets of N in G form a group G/N under the operation $(aN)(bN) = abN$. This group is called the **factor** or **quotient group** of G and N. Our first task is to prove that G/N is indeed a group.

Theorem 10.4. *Let N be a normal subgroup of a group G. The cosets of N in G form a group G/N of order $[G : N]$.*

PROOF.   The group operation on G/N is $(aN)(bN) = abN$. This operation must be shown to be well-defined; that is, group multiplication must be independent of the choice of coset representative. Let $aN = bN$ and $cN = dN$. We must show that

$$(aN)(cN) = acN = bdN = (bN)(dN).$$

Then $a = bn_1$ and $c = dn_2$ for some n_1 and n_2 in N. Hence,

$$\begin{aligned} acN &= bn_1 dn_2 N \\ &= bn_1 dN \\ &= bn_1 Nd \\ &= bNd \\ &= bdN. \end{aligned}$$

The remainder of the theorem is easy: $eN = N$ is the identity and $g^{-1}N$ is the inverse of gN. The order of G/N is, of course, the number of cosets of N in G. $\qquad\square$

It is very important to remember that the elements in a factor group are *sets of elements* in the original group.

Example 10.5. Consider the normal subgroup of S_3, $N = \{(1), (123), (132)\}$. The cosets of N in S_3 are N and $(12)N$. The factor group S_3/N has the following multiplication table.

	N	$(12)N$
N	N	$(12)N$
$(12)N$	$(12)N$	N

This group is isomorphic to $\mathbb{Z}_2$. At first, multiplying cosets seems both complicated and strange; however, notice that S_3/N is a smaller group. The factor group displays a certain amount of information about S_3. Actually, $N = A_3$, the group of even permutations, and $(12)N = \{(12), (13), (23)\}$ is the set of odd permutations. The information captured in G/N is parity; that is, multiplying two even or two odd permutations results in an even permutation, whereas multiplying an odd permutation by an even permutation yields an odd permutation.

Example 10.6. Consider the normal subgroup $3\mathbb{Z}$ of $\mathbb{Z}$. The cosets of $3\mathbb{Z}$ in $\mathbb{Z}$ are

$$0 + 3\mathbb{Z} = \{\ldots, -3, 0, 3, 6, \ldots\}$$
$$1 + 3\mathbb{Z} = \{\ldots, -2, 1, 4, 7, \ldots\}$$
$$2 + 3\mathbb{Z} = \{\ldots, -1, 2, 5, 8, \ldots\}.$$

The group $\mathbb{Z}/3\mathbb{Z}$ is given by the multiplication table below.

$+$	$0 + 3\mathbb{Z}$	$1 + 3\mathbb{Z}$	$2 + 3\mathbb{Z}$
$0 + 3\mathbb{Z}$	$0 + 3\mathbb{Z}$	$1 + 3\mathbb{Z}$	$2 + 3\mathbb{Z}$
$1 + 3\mathbb{Z}$	$1 + 3\mathbb{Z}$	$2 + 3\mathbb{Z}$	$0 + 3\mathbb{Z}$
$2 + 3\mathbb{Z}$	$2 + 3\mathbb{Z}$	$0 + 3\mathbb{Z}$	$1 + 3\mathbb{Z}$

In general, the subgroup $n\mathbb{Z}$ of $\mathbb{Z}$ is normal. The cosets of $\mathbb{Z}/n\mathbb{Z}$ are

$$n\mathbb{Z}$$
$$1 + n\mathbb{Z}$$
$$2 + n\mathbb{Z}$$
$$\vdots$$
$$(n-1) + n\mathbb{Z}.$$

The sum of the cosets $k + \mathbb{Z}$ and $l + \mathbb{Z}$ is $k + l + \mathbb{Z}$. Notice that we have written our cosets additively, because the group operation is integer addition.

Example 10.7. Consider the dihedral group D_n, generated by the two elements r and s, satisfying the relations

$$r^n = \mathrm{id}$$
$$s^2 = \mathrm{id}$$
$$srs = r^{-1}.$$

The element r actually generates the cyclic subgroup of rotations, R_n, of D_n. Since $srs^{-1} = srs = r^{-1} \in R_n$, the group of rotations is a normal subgroup of D_n; therefore, D_n/R_n is a group. Since there are exactly two elements in this group, it must be isomorphic to $\mathbb{Z}_2$.

10.2   The Simplicity of the Alternating Group

Of special interest are groups with no nontrivial normal subgroups. Such groups are called **simple groups**. Of course, we already have a whole class of examples of simple groups, $\mathbb{Z}_p$, where p is prime. These groups are trivially simple since they have no proper subgroups other than the subgroup consisting solely of the identity. Other examples of simple groups are not so easily found. We can, however, show that the alternating group, A_n, is simple for $n \geq 5$. The proof of this result requires several lemmas.

Lemma 10.8. *The alternating group A_n is generated by 3-cycles for $n \geq 3$.*

PROOF.   To show that the 3-cycles generate A_n, we need only show that any pair of transpositions can be written as the product of 3-cycles. Since $(ab) = (ba)$, every pair of transpositions must be one of the following:

$$(ab)(ab) = \mathrm{id}$$

$$(ab)(cd) = (acb)(acd)$$
$$(ab)(ac) = (acb). \qquad \square$$

Lemma 10.9. *Let N be a normal subgroup of A_n, where $n \geq 3$. If N contains a 3-cycle, then $N = A_n$.*

PROOF. We will first show that A_n is generated by 3-cycles of the specific form (ijk), where i and j are fixed in $\{1, 2, \ldots, n\}$ and we let k vary. Every 3-cycle is the product of 3-cycles of this form, since

$$(iaj) = (ija)^2$$
$$(iab) = (ijb)(ija)^2$$
$$(jab) = (ijb)^2(ija)$$
$$(abc) = (ija)^2(ijc)(ijb)^2(ija).$$

Now suppose that N is a nontrivial normal subgroup of A_n for $n \geq 3$ such that N contains a 3-cycle of the form (ija). Using the normality of N, we see that

$$[(ij)(ak)](ija)^2[(ij)(ak)]^{-1} = (ijk)$$

is in N. Hence, N must contain all of the 3-cycles (ijk) for $1 \leq k \leq n$. By Lemma 10.8, these 3-cycles generate A_n; hence, $N = A_n$. $\qquad \square$

Lemma 10.10. *For $n \geq 5$, every nontrivial normal subgroup N of A_n contains a 3-cycle.*

PROOF. Let σ be an arbitrary element in a normal subgroup N. There are several possible cycle structures for σ.

- σ is a 3-cycle.

- σ is the product of disjoint cycles, $\sigma = \tau(a_1 a_2 \cdots a_r) \in N$, where $r > 3$.

- σ is the product of disjoint cycles, $\sigma = \tau(a_1 a_2 a_3)(a_4 a_5 a_6)$.

- $\sigma = \tau(a_1 a_2 a_3)$, where τ is the product of disjoint 2-cycles.

- $\sigma = \tau(a_1 a_2)(a_3 a_4)$, where τ is the product of an even number of disjoint 2-cycles.

If σ is a 3-cycle, then we are done. If N contains a product of disjoint cycles, σ, and at least one of these cycles has length greater than 3, say $\sigma = \tau(a_1 a_2 \cdots a_r)$, then

$$(a_1 a_2 a_3)\sigma(a_1 a_2 a_3)^{-1}$$

is in N since N is normal; hence,

$$\sigma^{-1}(a_1 a_2 a_3)\sigma(a_1 a_2 a_3)^{-1}$$

is also in N. Since

$$\sigma^{-1}(a_1 a_2 a_3)\sigma(a_1 a_2 a_3)^{-1} = \sigma^{-1}(a_1 a_2 a_3)\sigma(a_1 a_3 a_2)$$
$$= (a_1 a_2 \cdots a_r)^{-1}\tau^{-1}(a_1 a_2 a_3)\tau(a_1 a_2 \cdots a_r)(a_1 a_3 a_2)$$
$$= (a_1 a_r a_{r-1} \cdots a_2)(a_1 a_2 a_3)(a_1 a_2 \cdots a_r)(a_1 a_3 a_2)$$
$$= (a_1 a_3 a_r),$$

N must contain a 3-cycle; hence, $N = A_n$.

Now suppose that N contains a disjoint product of the form

$$\sigma = \tau(a_1 a_2 a_3)(a_4 a_5 a_6).$$

Then

$$\sigma^{-1}(a_1 a_2 a_4)\sigma(a_1 a_2 a_4)^{-1} \in N$$

since

$$(a_1 a_2 a_4)\sigma(a_1 a_2 a_4)^{-1} \in N.$$

So

$$
\begin{aligned}
\sigma^{-1}(a_1 a_2 a_4)\sigma(a_1 a_2 a_4)^{-1} &= [\tau(a_1 a_2 a_3)(a_4 a_5 a_6)]^{-1}(a_1 a_2 a_4)\tau(a_1 a_2 a_3)(a_4 a_5 a_6)(a_1 a_2 a_4)^{-1} \\
&= (a_4 a_6 a_5)(a_1 a_3 a_2)\tau^{-1}(a_1 a_2 a_4)\tau(a_1 a_2 a_3)(a_4 a_5 a_6)(a_1 a_4 a_2) \\
&= (a_4 a_6 a_5)(a_1 a_3 a_2)(a_1 a_2 a_4)(a_1 a_2 a_3)(a_4 a_5 a_6)(a_1 a_4 a_2) \\
&= (a_1 a_4 a_2 a_6 a_3).
\end{aligned}
$$

So N contains a disjoint cycle of length greater than 3, and we can apply the previous case.

Suppose N contains a disjoint product of the form $\sigma = \tau(a_1 a_2 a_3)$, where τ is the product of disjoint 2-cycles. Since $\sigma \in N$, $\sigma^2 \in N$, and

$$
\begin{aligned}
\sigma^2 &= \tau(a_1 a_2 a_3)\tau(a_1 a_2 a_3) \\
&= (a_1 a_3 a_2).
\end{aligned}
$$

So N contains a 3-cycle.

The only remaining possible case is a disjoint product of the form

$$\sigma = \tau(a_1 a_2)(a_3 a_4),$$

where τ is the product of an even number of disjoint 2-cycles. But

$$\sigma^{-1}(a_1 a_2 a_3)\sigma(a_1 a_2 a_3)^{-1}$$

is in N since $(a_1 a_2 a_3)\sigma(a_1 a_2 a_3)^{-1}$ is in N; and so

$$
\begin{aligned}
\sigma^{-1}(a_1 a_2 a_3)\sigma(a_1 a_2 a_3)^{-1} &= \tau^{-1}(a_1 a_2)(a_3 a_4)(a_1 a_2 a_3)\tau(a_1 a_2)(a_3 a_4)(a_1 a_2 a_3)^{-1} \\
&= (a_1 a_3)(a_2 a_4).
\end{aligned}
$$

Since $n \geq 5$, we can find $b \in \{1, 2, \ldots, n\}$ such that $b \neq a_1, a_2, a_3, a_4$. Let $\mu = (a_1 a_3 b)$. Then

$$\mu^{-1}(a_1 a_3)(a_2 a_4)\mu(a_1 a_3)(a_2 a_4) \in N$$

and

$$
\begin{aligned}
\mu^{-1}(a_1 a_3)(a_2 a_4)\mu(a_1 a_3)(a_2 a_4) &= (a_1 b a_3)(a_1 a_3)(a_2 a_4)(a_1 a_3 b)(a_1 a_3)(a_2 a_4) \\
&= (a_1 a_3 b).
\end{aligned}
$$

Therefore, N contains a 3-cycle. This completes the proof of the lemma. $\qquad\square$

Theorem 10.11. *The alternating group, A_n, is simple for $n \geq 5$.*

PROOF. Let N be a normal subgroup of A_n. By Lemma 10.10, N contains a 3-cycle. By Lemma 10.9, $N = A_n$; therefore, A_n contains no proper nontrivial normal subgroups for $n \geq 5$. $\qquad\square$

One of the foremost problems of group theory has been to classify all simple finite groups. This problem is over a century old and has been solved only in the last few decades of the twentieth century. In a sense, finite simple groups are the building blocks of all finite groups. The first nonabelian simple groups to be discovered were the alternating groups. Galois was the first to prove that A_5 was simple. Later, mathematicians such as C. Jordan and L. E. Dickson found several infinite families of matrix groups that were simple. Other families of simple groups were discovered in the 1950s. At the turn of the century, William Burnside conjectured that all nonabelian simple groups must have even order. In 1963, W. Feit and J. Thompson proved Burnside's conjecture and published their results in the paper "Solvability of Groups of Odd Order," which appeared in the *Pacific Journal of Mathematics*. Their proof, running over 250 pages, gave impetus to a program in the 1960s and 1970s to classify all finite simple groups. Daniel Gorenstein was the organizer of this remarkable effort. One of the last simple groups was the "Monster," discovered by R. Greiss. The Monster, a 196,833 × 196,833 matrix group, is one of the 26 sporadic, or special, simple groups. These sporadic simple groups are groups that fit into no infinite family of simple groups. Some of the sporadic groups play an important role in physics.

10.3 Exercises

1.   For each of the following groups G, determine whether H is a normal subgroup of G. If H is a normal subgroup, write out a Cayley table for the factor group G/H.

(a) $G = S_4$ and $H = A_4$

(b) $G = A_5$ and $H = \{(1), (123), (132)\}$

(c) $G = S_4$ and $H = D_4$

(d) $G = Q_8$ and $H = \{1, -1, I, -I\}$

(e) $G = \mathbb{Z}$ and $H = 5\mathbb{Z}$

2.   Find all the subgroups of D_4. Which subgroups are normal? What are all the factor groups of D_4 up to isomorphism?

3.   Find all the subgroups of the quaternion group, Q_8. Which subgroups are normal? What are all the factor groups of Q_8 up to isomorphism?

4.   Let T be the group of nonsingular upper triangular 2×2 matrices with entries in $\mathbb{R}$; that is, matrices of the form

$$\begin{pmatrix} a & b \\ 0 & c \end{pmatrix},$$

where a, b, $c \in \mathbb{R}$ and $ac \neq 0$. Let U consist of matrices of the form

$$\begin{pmatrix} 1 & x \\ 0 & 1 \end{pmatrix},$$

where $x \in \mathbb{R}$.

(a) Show that U is a subgroup of T.

(b) Prove that U is abelian.

(c) Prove that U is normal in T.

(d) Show that T/U is abelian.

(e) Is T normal in $GL_2(\mathbb{R})$?

5. Show that the intersection of two normal subgroups is a normal subgroup.

6. If G is abelian, prove that G/H must also be abelian.

7. Prove or disprove: If H is a normal subgroup of G such that H and G/H are abelian, then G is abelian.

8. If G is cyclic, prove that G/H must also be cyclic.

9. Prove or disprove: If H and G/H are cyclic, then G is cyclic.

10. Let H be a subgroup of index 2 of a group G. Prove that H must be a normal subgroup of G. Conclude that S_n is not simple for $n \geq 3$.

11. If a group G has exactly one subgroup H of order k, prove that H is normal in G.

12. Define the **centralizer** of an element g in a group G to be the set

$$C(g) = \{x \in G : xg = gx\}.$$

Show that $C(g)$ is a subgroup of G. If g generates a normal subgroup of G, prove that $C(g)$ is normal in G.

13. Recall that the **center** of a group G is the set

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

(a) Calculate the center of S_3.

(b) Calculate the center of $GL_2(\mathbb{R})$.

(c) Show that the center of any group G is a normal subgroup of G.

(d) If $G/Z(G)$ is cyclic, show that G is abelian.

14. Let G be a group and let $G' = \langle aba^{-1}b^{-1} \rangle$; that is, G' is the subgroup of all finite products of elements in G of the form $aba^{-1}b^{-1}$. The subgroup G' is called the **commutator subgroup** of G.

(a) Show that G' is a normal subgroup of G.

(b) Let N be a normal subgroup of G. Prove that G/N is abelian if and only if N contains the commutator subgroup of G.

10.4 Sage

Sage has several convenient functions that will allow us to investigate quickly if a subgroup is normal, and if so, the nature of the resulting quotient group. But for an initial understanding, we can also work with the raw cosets. Let us get our hands dirty first, then learn about the easy way.

Multiplying Cosets

The definiton of a factor group requires a normal subgroup, and then we *define* a way to "multiply" two cosets of the subgroup to produce another coset. It is important to realize that we can interpret the definition of a normal subgroup to be *exactly* the condition we need for our new multiplication to be workable. We will do two examples — first with a normal subgroup, then with a subgroup that is not normal.

Consider the dihedral group D_8 that is the symmetry group of an 8-gon. If we take the element that creates a quarter-turn, we can use it generate a cyclic subgroup of order 4. This will be a normal subgroup (trust us for the moment on this). First, build the (right) cosets (notice there is no output):

```
G = DihedralGroup(8)
quarter_turn = G('(1,3,5,7)(2,4,6,8)')
S = G.subgroup([quarter_turn])
C = G.cosets(S)
```

So C is a list of lists, with every element of the group G occuring exactly once somewhere. You could ask Sage to print out C for you if you like, but we will try to avoid that here. We want to multiply two cosets (lists) together. How do we do this? Take *any* element out of the first list, and *any* element out of the second list and multiply them together (which we know how to do since they are elements of G). Now we have an element of G. What do we do with this element, since we really want a coset as the result of the product of two cosets? Simple — we see which coset the product is in. Let us give it a try. We will multiply coset 1 with coset 3 (there are 4 cosets by Lagrange's Theorem). Study the following code carefully to see if you can understand what it is doing, and *then* read the explanation that follows.

```
p = C[1][0]*C[3][0]
[i for i in srange(len(C)) if p in C[i]]
```

```
[2]
```

What have we accomplished? In the first line we create p as the product of two group elements, one from coset 1 and one from coset 3 (C[1], C[3]). Since we can choose *any* element from each coset, we choose the first element of each (C[ ][0]). Then we count our way through all the cosets, selecting only cosets that contain p. Since p will only be in one coset, we expect a list with just one element. Here, our one-element list contains only 2. So we say the product of coset 1 and coset 3 is coset 2.

The point here is that this result (coset 1 times coset 3 is coset 2) should always be the same, *no matter which elements we pick from the two cosets to form* p. So let us do it again, but this time we will not simply choose the first element from each of coset 1 and coset 3, instead we will choose the third element of coset 1 and the second element of coset 3 (remember, we are counting from zero!).

```
p = C[1][2]*C[3][1]
[i for i in srange(len(C)) if p in C[i]]
```

```
[2]
```

Good. We have the same result. If you are still trusting us on S being a normal subgroup of G, then this is the result that the theory predicts. Make a copy of the above compute cell and try other choices for the representatives of each coset. Then try the product of other cosets, with varying representatives.Now is a good time to introduce a way to extend Sage and add new functions. We will design a coset-multiplication function. Read the following carefully and then see the subsequent explanation.

```
def coset_product(i, j, C):
    p = C[i][0]*C[j][0]
    c = [k for k in srange(len(C)) if p in C[k]]
    return c[0]
```

The first line creates a new Sage function named `coset_product`. This is accomplished with the word `def`, and note the colon ending the line. The inputs to the function are the numbers of the cosets we want to multiply and the complete list of the cosets. The middle two lines should look familiar from above. We know `c` is a one-element list, so `c[0]` will extract this one coset number, and `return` is what determines that this is the output of the function. Notice that the indentation above must be exactly as shown. We could have written all this computation on a single line without making a new function, but that begins to get unwieldly. You need to execute the code block above to actually *define* the function, and there will be no output if successful. Now we can use our new function to repeat our work above:

```
coset_product(1, 3, C)
```

2

Now you know the basics of how to add onto Sage and do much more than it was designed for. And with some practice, you could suggest and contribute new functions to Sage, since it is an open source project. Nice.

Now let us examine a situation where the subgroup is not normal. So we will see that our definition of coset multiplication is insufficient in this case. And realize that our new `coset_product` function is also useless since it assumes the cosets come from a normal subgroup.

Consider the alternating group A_4 which we can interpet as the symmetry group of a tetrahedron. For a subgroup, take an element that fixes one vertex and rotates the opposite face — this will generate a cyclic subgroup of order 3, and by Lagrange's Theorem we will get four cosets. We compute them here. (Again, no output is requested.)

```
G = AlternatingGroup(4)
face_turn = G("(1,2,3)")
S = G.subgroup([face_turn])
C = G.cosets(S)
```

Again, let's consider the product of coset 1 and coset 3:

```
p = C[1][0]*C[3][0]
[i for i in srange(len(C)) if p in C[i]]
```

[0]

Again, but now for coset 3, choose the second element of the coset to produce the product p:

```
p = C[1][0]*C[3][1]
[i for i in srange(len(C)) if p in C[i]]
```

[2]

So, is the product of coset 1 and coset 3 equal to coset 0 or coset 2? We cannot say! So there is *no way* to construct a quotient group for this subgroup. You can experiment some more with this subgroup, but in some sense, we are done with this example — there is nothing left to say.

Sage Methods for Normal Subgroups

You can easily ask Sage if a subgroup is normal or not. This is viewed as a property of the subgroup, but you must tell Sage what the "supergroup" is, since the answer can change depending on this value. (For example H.is_normal(H) will always be True.) Here are our two examples from above.

```
G = DihedralGroup(8)
quarter_turn = G('(1,3,5,7)(2,4,6,8)')
S = G.subgroup([quarter_turn])
S.is_normal(G)
```

True

```
G = AlternatingGroup(4)
face_turn = G("(1,2,3)")
S = G.subgroup([face_turn])
S.is_normal(G)
```

False

The text proves in Section 10.2 that A_5 is simple, i.e. A_5 has no normal subgroups. We could build every subgroup of A_5 and ask if it is normal in A_5 using the .is_normal() method. But Sage has this covered for us already.

```
G = AlternatingGroup(5)
G.is_simple()
```

True

We can also build a quotient group when we have a normal subgroup.

```
G = DihedralGroup(8)
quarter_turn = G('(1,3,5,7)(2,4,6,8)')
S = G.subgroup([quarter_turn])
Q = G.quotient(S)
Q
```

Permutation Group with generators [(1,2)(3,4), (1,3)(2,4)]

This is useful, but also a bit unsettling. We have the quotient group, but any notion of cosets has been lost, since Q is returned as a new permutation group on a different set of symbols. We cannot presume that the numbers used for the new permutation group Q bear any resemblance to the cosets we get from the .cosets() method. But we can see that the quotient group is described as a group generated by two elements of order two. We could ask for the order of the group, or by Lagrange's Theorem we know the quotient has order 4. We can say now that there are only two groups of order four, the cyclic group of order 4 and a non-cyclic group of order 4, known to us as the Klein 4-group or $\mathbb{Z}_2 \times \mathbb{Z}_2$. This quotient group looks like the non-cyclic one since the cyclic group of order 4 has just one element of order 2. Let us see what Sage says.

```
Q.is_isomorphic(KleinFourGroup())
```

True

Yes, that's it.

Finally, Sage can build us a list of all of the normal subgroups of a group. The list of groups themselves, as we have seen before, is sometimes an overwhelming amount of information. We will demonstrate by just listing the orders of the normal subgroups produced.

```
G = DihedralGroup(8)
N = G.normal_subgroups()
[H.order() for H in N]
```

```
[1, 2, 4, 8, 8, 8, 16]
```

So, in particular, we see that our "quarter-turn" subgroup is the *only* normal subgroup of order 4 in this group.

10.5   Sage Exercises

1.   Build every subgroup of the alternating group on 5 symbols, A_5, and check that each is not a normal subgroup (except for the two trivial cases). This command might take a couple seconds to run. Compare this with the time needed to run the `.is_simple()` method and realize that there is a significant amount of theory and cleverness brought to bear in speeding up commands like this. (It is possible that your Sage installation lacks GAP's "Table of Marks" library and you will be unable to compute the list of subgroups.)

2.   Consider the quotient group of the group of symmetries of an 8-gon, formed with the cyclic subgroup of order 4 generated by a quarter-turn. Use the `coset_product` function to determine the Cayley table for this quotient group. Use the number of each coset, as produced by the `.cosets()` method as names for the elements of the quotient group. You will need to build the table "by hand" as there is no easy way to have Sage's Cayley table command do this one for you. You can build a table in the Sage Notebook pop-up editor (shift-click on a blue line) or you might read the documentation of the `html.table()` method.

3.   Consider the cyclic subgroup of order 4 in the symmetries of an 8-gon. Verify that the subgroup is normal by first building the raw left and right cosets (without using the `.cosets()` method) and then checking their equality in Sage, all with a single command that employs sorting with the `sorted()` command.

4.   Again, use the same cyclic subgroup of order 4 in the group of symmetries of an 8-gon. Check that the subgroup is normal by using part (2) of Theorem 10.3. Construct a one-line command that does the complete check and returns `True`. Maybe sort the elements of the subgroup `S` first, then slowly build up the necessary lists, commands, and conditions in steps. Notice that this check does not require ever building the cosets.

5.   Repeat the demonstration from the previous subsection that for the symmetries of a tetrahedron, a cyclic subgroup of order 3 results in an undefined coset multiplication. Above, the default setting for the `.cosets()` method builds right cosets — but in this problem, work instead with left cosets. You need to choose two cosets to multiply, and then demonstrate two choices for representatives that lead to different results for the product of the cosets.

6.   Construct some dihedral groups of order $2n$ (i.e. symmetries of an n-gon, D_n in the text, `DihedralGroup(n)` in Sage). Maybe all of them for $3 \leq n \leq 100$. For each dihedral group, construct a list of the orders of each of the normal subgroups (so use `.normal_subgroups()`). You may need to wait ten or twenty seconds for this to finish - be patient. Observe enough examples to hypothesize a pattern to your observations, check your hypothesis against each of your examples and then state your hypothesis clearly.

Can you predict how many normal subgroups there are in the dihedral group D_{470448} without using Sage to build all the normal subgroups? Can you *describe* all of the normal subgroups of a dihedral group in a way that would let us predict all of the normal subgroups of D_{470448} without using Sage?

11

Homomorphisms

One of the basic ideas of algebra is the concept of a homomorphism, a natural generalization of an isomorphism. If we relax the requirement that an isomorphism of groups be bijective, we have a homomorphism.

11.1 Group Homomorphisms

A **homomorphism** between groups $(G, \cdot)$ and $(H, \circ)$ is a map $\phi : G \to H$ such that

$$\phi(g_1 \cdot g_2) = \phi(g_1) \circ \phi(g_2)$$

for $g_1, g_2 \in G$. The range of ϕ in H is called the **homomorphic image** of ϕ.

Two groups are related in the strongest possible way if they are isomorphic; however, a weaker relationship may exist between two groups. For example, the symmetric group S_n and the group $\mathbb{Z}_2$ are related by the fact that S_n can be divided into even and odd permutations that exhibit a group structure like that $\mathbb{Z}_2$, as shown in the following multiplication table.

	even	odd
even	even	odd
odd	odd	even

We use homomorphisms to study relationships such as the one we have just described.

Example 11.1. Let G be a group and $g \in G$. Define a map $\phi : \mathbb{Z} \to G$ by $\phi(n) = g^n$. Then ϕ is a group homomorphism, since

$$\phi(m + n) = g^{m+n} = g^m g^n = \phi(m)\phi(n).$$

This homomorphism maps $\mathbb{Z}$ onto the cyclic subgroup of G generated by g.

Example 11.2. Let $G = GL_2(\mathbb{R})$. If

$$A = \begin{pmatrix} a & b \\ c & d \end{pmatrix}$$

is in G, then the determinant is nonzero; that is, $\det(A) = ad - bc \neq 0$. Also, for any two elements A and B in G, $\det(AB) = \det(A)\det(B)$. Using the determinant, we can define a homomorphism $\phi : GL_2(\mathbb{R}) \to \mathbb{R}^*$ by $A \mapsto \det(A)$.

Example 11.3. Recall that the circle group $\mathbb{T}$ consists of all complex numbers z such that $|z| = 1$. We can define a homomorphism ϕ from the additive group of real numbers $\mathbb{R}$ to $\mathbb{T}$ by $\phi : \theta \mapsto \cos\theta + i\sin\theta$. Indeed,

$$
\begin{aligned}
\phi(\alpha + \beta) &= \cos(\alpha + \beta) + i\sin(\alpha + \beta) \\
&= (\cos\alpha\cos\beta - \sin\alpha\sin\beta) + i(\sin\alpha\cos\beta + \cos\alpha\sin\beta) \\
&= (\cos\alpha + i\sin\alpha)(\cos\beta + i\sin\beta) \\
&= \phi(\alpha)\phi(\beta).
\end{aligned}
$$

Geometrically, we are simply wrapping the real line around the circle in a group-theoretic fashion.

The following proposition lists some basic properties of group homomorphisms.

Proposition 11.4. *Let $\phi : G_1 \to G_2$ be a homomorphism of groups. Then*

1. *If e is the identity of G_1, then $\phi(e)$ is the identity of G_2;*

2. *For any element $g \in G_1$, $\phi(g^{-1}) = [\phi(g)]^{-1}$;*

3. *If H_1 is a subgroup of G_1, then $\phi(H_1)$ is a subgroup of G_2;*

4. *If H_2 is a subgroup of G_2, then $\phi^{-1}(H_2) = \{g \in G_1 : \phi(g) \in H_2\}$ is a subgroup of G_1. Furthermore, if H_2 is normal in G_2, then $\phi^{-1}(H_2)$ is normal in G_1.*

PROOF.   (1) Suppose that e and e' are the identities of G_1 and G_2, respectively; then

$$
e'\phi(e) = \phi(e) = \phi(ee) = \phi(e)\phi(e).
$$

By cancellation, $\phi(e) = e'$.

(2) This statement follows from the fact that

$$
\phi(g^{-1})\phi(g) = \phi(g^{-1}g) = \phi(e) = e'.
$$

(3) The set $\phi(H_1)$ is nonempty since the identity of G_2 is in $\phi(H_1)$. Suppose that H_1 is a subgroup of G_1 and let x and y be in $\phi(H_1)$. There exist elements $a, b \in H_1$ such that $\phi(a) = x$ and $\phi(b) = y$. Since

$$
xy^{-1} = \phi(a)[\phi(b)]^{-1} = \phi(ab^{-1}) \in \phi(H_1),
$$

$\phi(H_1)$ is a subgroup of G_2 by Proposition 3.31.

(4) Let H_2 be a subgroup of G_2 and define H_1 to be $\phi^{-1}(H_2)$; that is, H_1 is the set of all $g \in G_1$ such that $\phi(g) \in H_2$. The identity is in H_1 since $\phi(e) = e'$. If a and b are in H_1, then $\phi(ab^{-1}) = \phi(a)[\phi(b)]^{-1}$ is in H_2 since H_2 is a subgroup of G_2. Therefore, $ab^{-1} \in H_1$ and H_1 is a subgroup of G_1. If H_2 is normal in G_2, we must show that $g^{-1}hg \in H_1$ for $h \in H_1$ and $g \in G_1$. But

$$
\phi(g^{-1}hg) = [\phi(g)]^{-1}\phi(h)\phi(g) \in H_2,
$$

since H_2 is a normal subgroup of G_2. Therefore, $g^{-1}hg \in H_1$. $\square$

Let $\phi : G \to H$ be a group homomorphism and suppose that e is the identity of H. By Proposition 11.4, $\phi^{-1}(\{e\})$ is a subgroup of G. This subgroup is called the ***kernel*** of ϕ and will be denoted by $\ker\phi$. In fact, this subgroup is a normal subgroup of G since the trivial subgroup is normal in H. We state this result in the following theorem, which says that with every homomorphism of groups we can naturally associate a normal subgroup.

Theorem 11.5. *Let $\phi : G \to H$ be a group homomorphism. Then the kernel of ϕ is a normal subgroup of G.*

Example 11.6. Let us examine the homomorphism $\phi : GL_2(\mathbb{R}) \to \mathbb{R}^*$ defined by $A \mapsto \det(A)$. Since 1 is the identity of $\mathbb{R}^*$, the kernel of this homomorphism is all 2×2 matrices having determinant one. That is, $\ker \phi = SL_2(\mathbb{R})$.

Example 11.7. The kernel of the group homomorphism $\phi : \mathbb{R} \to \mathbb{C}^*$ defined by $\phi(\theta) = \cos \theta + i \sin \theta$ is $\{2\pi n : n \in \mathbb{Z}\}$. Notice that $\ker \phi \cong \mathbb{Z}$.

Example 11.8. Suppose that we wish to determine all possible homomorphisms ϕ from $\mathbb{Z}_7$ to $\mathbb{Z}_{12}$. Since the kernel of ϕ must be a subgroup of $\mathbb{Z}_7$, there are only two possible kernels, $\{0\}$ and all of $\mathbb{Z}_7$. The image of a subgroup of $\mathbb{Z}_7$ must be a subgroup of $\mathbb{Z}_{12}$. Hence, there is no injective homomorphism; otherwise, $\mathbb{Z}_{12}$ would have a subgroup of order 7, which is impossible. Consequently, the only possible homomorphism from $\mathbb{Z}_7$ to $\mathbb{Z}_{12}$ is the one mapping all elements to zero.

Example 11.9. Let G be a group. Suppose that $g \in G$ and ϕ is the homomorphism from $\mathbb{Z}$ to G given by $\phi(n) = g^n$. If the order of g is infinite, then the kernel of this homomorphism is $\{0\}$ since ϕ maps $\mathbb{Z}$ onto the cyclic subgroup of G generated by g. However, if the order of g is finite, say n, then the kernel of ϕ is $n\mathbb{Z}$.

11.2 The Isomorphism Theorems

Although it is not evident at first, factor groups correspond exactly to homomorphic images, and we can use factor groups to study homomorphisms. We already know that with every group homomorphism $\phi : G \to H$ we can associate a normal subgroup of G, $\ker \phi$. The converse is also true; that is, every normal subgroup of a group G gives rise to homomorphism of groups.

Let H be a normal subgroup of G. Define the ***natural*** or ***canonical homomorphism***

$$\phi : G \to G/H$$

by

$$\phi(g) = gH.$$

This is indeed a homomorphism, since

$$\phi(g_1 g_2) = g_1 g_2 H = g_1 H g_2 H = \phi(g_1)\phi(g_2).$$

The kernel of this homomorphism is H. The following theorems describe the relationships between group homomorphisms, normal subgroups, and factor groups.

Theorem 11.10 First Isomorphism Theorem. *If $\psi : G \to H$ is a group homomorphism with $K = \ker \psi$, then K is normal in G. Let $\phi : G \to G/K$ be the canonical homomorphism. Then there exists a unique isomorphism $\eta : G/K \to \psi(G)$ such that $\psi = \eta\phi$.*

PROOF. We already know that K is normal in G. Define $\eta : G/K \to \psi(G)$ by $\eta(gK) = \psi(g)$. We first show that η is a well-defined map. If $g_1 K = g_2 K$, then for some $k \in K$, $g_1 k = g_2$; consequently,

$$\eta(g_1 K) = \psi(g_1) = \psi(g_1)\psi(k) = \psi(g_1 k) = \psi(g_2) = \eta(g_2 K).$$

Thus, η does not depend on the choice of coset representatives and the map $\eta : G/K \to \psi(G)$ is uniquely defined since $\psi = \eta\phi$. We must also show that η is a homomorphism, but

$$\begin{aligned}
\eta(g_1 K g_2 K) &= \eta(g_1 g_2 K) \\
&= \psi(g_1 g_2) \\
&= \psi(g_1)\psi(g_2) \\
&= \eta(g_1 K)\eta(g_2 K).
\end{aligned}$$

Clearly, η is onto $\psi(G)$. To show that η is one-to-one, suppose that $\eta(g_1 K) = \eta(g_2 K)$. Then $\psi(g_1) = \psi(g_2)$. This implies that $\psi(g_1^{-1} g_2) = e$, or $g_1^{-1} g_2$ is in the kernel of ψ; hence, $g_1^{-1} g_2 K = K$; that is, $g_1 K = g_2 K$. $\qquad\qquad\square$

Mathematicians often use diagrams called **commutative diagrams** to describe such theorems. The following diagram "commutes" since $\psi = \eta\phi$.

$$
\begin{array}{ccc}
G & \xrightarrow{\ \ \psi\ \ } & H \\
 & \phi \searrow \ \ \nearrow \eta & \\
 & G/K &
\end{array}
$$

Example 11.11. Let G be a cyclic group with generator g. Define a map $\phi : \mathbb{Z} \to G$ by $n \mapsto g^n$. This map is a surjective homomorphism since

$$\phi(m + n) = g^{m+n} = g^m g^n = \phi(m)\phi(n).$$

Clearly ϕ is onto. If $|g| = m$, then $g^m = e$. Hence, $\ker\phi = m\mathbb{Z}$ and $\mathbb{Z}/\ker\phi = \mathbb{Z}/m\mathbb{Z} \cong G$. On the other hand, if the order of g is infinite, then $\ker\phi = 0$ and ϕ is an isomorphism of G and $\mathbb{Z}$. Hence, two cyclic groups are isomorphic exactly when they have the same order. Up to isomorphism, the only cyclic groups are $\mathbb{Z}$ and $\mathbb{Z}_n$.

Theorem 11.12 Second Isomorphism Theorem. *Let H be a subgroup of a group G (not necessarily normal in G) and N a normal subgroup of G. Then HN is a subgroup of G, $H \cap N$ is a normal subgroup of H, and*

$$H/H \cap N \cong HN/N.$$

PROOF.   We will first show that $HN = \{hn : h \in H, n \in N\}$ is a subgroup of G. Suppose that $h_1 n_1, h_2 n_2 \in HN$. Since N is normal, $(h_2)^{-1} n_1 h_2 \in N$. So

$$(h_1 n_1)(h_2 n_2) = h_1 h_2((h_2)^{-1} n_1 h_2)n_2$$

is in HN. The inverse of $hn \in HN$ is in HN since

$$(hn)^{-1} = n^{-1} h^{-1} = h^{-1}(hn^{-1}h^{-1}).$$

Next, we prove that $H \cap N$ is normal in H. Let $h \in H$ and $n \in H \cap N$. Then $h^{-1}nh \in H$ since each element is in H. Also, $h^{-1}nh \in N$ since N is normal in G; therefore, $h^{-1}nh \in H \cap N$.

Now define a map ϕ from H to HN/N by $h \mapsto hN$. The map ϕ is onto, since any coset $hnN = hN$ is the image of h in H. We also know that ϕ is a homomorphism because

$$\phi(hh') = hh'N = hNh'N = \phi(h)\phi(h').$$

By the First Isomorphism Theorem, the image of ϕ is isomorphic to $H/\ker\phi$; that is,

$$HN/N = \phi(H) \cong H/\ker\phi.$$

Since

$$\ker\phi = \{h \in H : h \in N\} = H \cap N,$$

$HN/N = \phi(H) \cong H/H \cap N.$ □

Theorem 11.13 Correspondence Theorem. *Let N be a normal subgroup of a group G. Then $H \mapsto H/N$ is a one-to-one correspondence between the set of subgroups H containing N and the set of subgroups of G/N. Furthermore, the normal subgroups of G containing N correspond to normal subgroups of G/N.*

PROOF. Let H be a subgroup of G containing N. Since N is normal in H, H/N makes sense. Let aN and bN be elements of H/N. Then $(aN)(b^{-1}N) = ab^{-1}N \in H/N$; hence, H/N is a subgroup of G/N.

Let S be a subgroup of G/N. This subgroup is a set of cosets of N. If $H = \{g \in G : gN \in S\}$, then for $h_1, h_2 \in H$, we have that $(h_1N)(h_2N) = h_1h_2N \in S$ and $h_1^{-1}N \in S$. Therefore, H must be a subgroup of G. Clearly, H contains N. Therefore, $S = H/N$. Consequently, the map $H \mapsto H/N$ is onto.

Suppose that H_1 and H_2 are subgroups of G containing N such that $H_1/N = H_2/N$. If $h_1 \in H_1$, then $h_1N \in H_1/N$. Hence, $h_1N = h_2N \subset H_2$ for some h_2 in H_2. However, since N is contained in H_2, we know that $h_1 \in H_2$ or $H_1 \subset H_2$. Similarly, $H_2 \subset H_1$. Since $H_1 = H_2$, the map $H \mapsto H/N$ is one-to-one.

Suppose that H is normal in G and N is a subgroup of H. Then it is easy to verify that the map $G/N \to G/H$ defined by $gN \mapsto gH$ is a homomorphism. The kernel of this homomorphism is H/N, which proves that H/N is normal in G/N.

Conversely, suppose that H/N is normal in G/N. The homomorphism given by

$$G \to G/N \to \frac{G/N}{H/N}$$

has kernel H. Hence, H must be normal in G. □

Notice that in the course of the proof of Theorem 11.13, we have also proved the following theorem.

Theorem 11.14 Third Isomorphism Theorem. *Let G be a group and N and H be normal subgroups of G with $N \subset H$. Then*

$$G/H \cong \frac{G/N}{H/N}.$$

Example 11.15. By the Third Isomorphism Theorem,

$$\mathbb{Z}/m\mathbb{Z} \cong (\mathbb{Z}/mn\mathbb{Z})/(m\mathbb{Z}/mn\mathbb{Z}).$$

Since $|\mathbb{Z}/mn\mathbb{Z}| = mn$ and $|\mathbb{Z}/m\mathbb{Z}| = m$, we have $|m\mathbb{Z}/mn\mathbb{Z}| = n$.

11.3   Exercises

1.   Prove that $\det(AB) = \det(A)\det(B)$ for $A, B \in GL_2(\mathbb{R})$.   This shows that the determinant is a homomorphism from $GL_2(\mathbb{R})$ to $\mathbb{R}^*$.

2.   Which of the following maps are homomorphisms? If the map is a homomorphism, what is the kernel?

(a) $\phi : \mathbb{R}^* \to GL_2(\mathbb{R})$ defined by

$$\phi(a) = \begin{pmatrix} 1 & 0 \\ 0 & a \end{pmatrix}$$

(b) $\phi : \mathbb{R} \to GL_2(\mathbb{R})$ defined by

$$\phi(a) = \begin{pmatrix} 1 & 0 \\ a & 1 \end{pmatrix}$$

(c) $\phi : GL_2(\mathbb{R}) \to \mathbb{R}$ defined by

$$\phi\left(\begin{pmatrix} a & b \\ c & d \end{pmatrix}\right) = a + d$$

(d) $\phi : GL_2(\mathbb{R}) \to \mathbb{R}^*$ defined by

$$\phi\left(\begin{pmatrix} a & b \\ c & d \end{pmatrix}\right) = ad - bc$$

(e) $\phi : \mathbb{M}_2(\mathbb{R}) \to \mathbb{R}$ defined by

$$\phi\left(\begin{pmatrix} a & b \\ c & d \end{pmatrix}\right) = b,$$

where $\mathbb{M}_2(\mathbb{R})$ is the additive group of 2×2 matrices with entries in $\mathbb{R}$.

3.   Let A be an $m \times n$ matrix. Show that matrix multiplication, $x \mapsto Ax$, defines a homomorphism $\phi : \mathbb{R}^n \to \mathbb{R}^m$.

4.   Let $\phi : \mathbb{Z} \to \mathbb{Z}$ be given by $\phi(n) = 7n$. Prove that ϕ is a group homomorphism. Find the kernel and the image of ϕ.

5.   Describe all of the homomorphisms from $\mathbb{Z}_{24}$ to $\mathbb{Z}_{18}$.

6.   Describe all of the homomorphisms from $\mathbb{Z}$ to $\mathbb{Z}_{12}$.

7.   In the group $\mathbb{Z}_{24}$, let $H = \langle 4 \rangle$ and $N = \langle 6 \rangle$.

(a) List the elements in HN (we usually write $H + N$ for these additive groups) and $H \cap N$.

(b) List the cosets in HN/N, showing the elements in each coset.

(c) List the cosets in $H/(H \cap N)$, showing the elements in each coset.

(d) Give the correspondence between HN/N and $H/(H \cap N)$ described in the proof of the Second Isomorphism Theorem.

8. If G is an abelian group and $n \in \mathbb{N}$, show that $\phi : G \to G$ defined by $g \mapsto g^n$ is a group homomorphism.

9. If $\phi : G \to H$ is a group homomorphism and G is abelian, prove that $\phi(G)$ is also abelian.

10. If $\phi : G \to H$ is a group homomorphism and G is cyclic, prove that $\phi(G)$ is also cyclic.

11. Show that a homomorphism defined on a cyclic group is completely determined by its action on the generator of the group.

12. If a group G has exactly one subgroup H of order k, prove that H is normal in G.

13. Prove or disprove: $\mathbb{Q}/\mathbb{Z} \cong \mathbb{Q}$.

14. Let G be a finite group and N a normal subgroup of G. If H is a subgroup of G/N, prove that $\phi^{-1}(H)$ is a subgroup in G of order $|H| \cdot |N|$, where $\phi : G \to G/N$ is the canonical homomorphism.

15. Let G_1 and G_2 be groups, and let H_1 and H_2 be normal subgroups of G_1 and G_2 respectively. Let $\phi : G_1 \to G_2$ be a homomorphism. Show that ϕ induces a natural homomorphism $\overline{\phi} : (G_1/H_1) \to (G_2/H_2)$ if $\phi(H_1) \subset H_2$.

16. If H and K are normal subgroups of G and $H \cap K = \{e\}$, prove that G is isomorphic to a subgroup of $G/H \times G/K$.

17. Let $\phi : G_1 \to G_2$ be a surjective group homomorphism. Let H_1 be a normal subgroup of G_1 and suppose that $\phi(H_1) = H_2$. Prove or disprove that $G_1/H_1 \cong G_2/H_2$.

18. Let $\phi : G \to H$ be a group homomorphism. Show that ϕ is one-to-one if and only if $\phi^{-1}(e) = \{e\}$.

19. Given a homomorphism $\phi : G \to H$ define a relation $\sim$ on G by $a \sim b$ if $\phi(a) = \phi(b)$ for $a, b \in G$. Show this relation is an equivalence relation and describe the equivalence classes.

11.4 Additional Exercises: Automorphisms

1. Let $\operatorname{Aut}(G)$ be the set of all automorphisms of G; that is, isomorphisms from G to itself. Prove this set forms a group and is a subgroup of the group of permutations of G; that is, $\operatorname{Aut}(G) \leq S_G$.

2. An ***inner automorphism*** of G,

$$i_g : G \to G,$$

is defined by the map

$$i_g(x) = gxg^{-1},$$

for $g \in G$. Show that $i_g \in \operatorname{Aut}(G)$.

3. The set of all inner automorphisms is denoted by $\operatorname{Inn}(G)$. Show that $\operatorname{Inn}(G)$ is a subgroup of $\operatorname{Aut}(G)$.

4. Find an automorphism of a group G that is not an inner automorphism.

5. Let G be a group and i_g be an inner automorphism of G, and define a map

$$G \to \operatorname{Aut}(G)$$

by

$$g \mapsto i_g.$$

Prove that this map is a homomorphism with image $\text{Inn}(G)$ and kernel $Z(G)$. Use this result to conclude that

$$G/Z(G) \cong \text{Inn}(G).$$

6.   Compute $\text{Aut}(S_3)$ and $\text{Inn}(S_3)$. Do the same thing for D_4.

7.   Find all of the homomorphisms $\phi : \mathbb{Z} \to \mathbb{Z}$. What is $\text{Aut}(\mathbb{Z})$?

8.   Find all of the automorphisms of $\mathbb{Z}_8$. Prove that $\text{Aut}(\mathbb{Z}_8) \cong U(8)$.

9.   For $k \in \mathbb{Z}_n$, define a map $\phi_k : \mathbb{Z}_n \to \mathbb{Z}_n$ by $a \mapsto ka$. Prove that ϕ_k is a homomorphism.

10.   Prove that ϕ_k is an isomorphism if and only if k is a generator of $\mathbb{Z}_n$.

11.   Show that every automorphism of $\mathbb{Z}_n$ is of the form ϕ_k, where k is a generator of $\mathbb{Z}_n$.

12.   Prove that $\psi : U(n) \to \text{Aut}(\mathbb{Z}_n)$ is an isomorphism, where $\psi : k \mapsto \phi_k$.

11.5 Sage

Sage is able to create homomorphisms (and by extension, isomorphisms and automorphisms) between finite permutation groups. There is a limited supply of commands then available to manipulate these functions, but we can still illustrate many of the ideas in this chapter.

Homomorphisms

The principal device for creating a homomorphism is to specify the specific images of the set of generators for the domain. Consider cyclic groups of order 12 and 20:

$$G = \{a^i | a^{12} = e\} \qquad\qquad H = \{x^i | x^{20} = e\}$$

and define a homomorphism by just defining the image of the generator of G, and define the rest of the mapping by extending the mapping via the operation-preserving property of a homomorphism.

$$\phi : G \to H, \quad \phi(a) = x^5$$
$$\Rightarrow \quad \phi(a^i) = \phi(a)^i = (x^5)^i = x^{5i}$$

The constructor `PermutationGroupMorphism` requires the two groups, then a list of images for each generator (in order!), and then will create the homomorphism. Note that we can then use the result as a function. In the example below, we first verify that `C12` has a single generator (no surprise there), which we then send to a particular element of order 4 in the codomain. Sage then constructs the unique homomorphism that is consistent with this requirement.

```
C12 = CyclicPermutationGroup(12)
C20 = CyclicPermutationGroup(20)
domain_gens = C12.gens()
[g.order() for g in domain_gens]
```

```
[12]
```

```
x = C20.gen(0)
y = x^5
y.order()
```

4

```
phi = PermutationGroupMorphism(C12, C20, [y])
phi
```

```
Permutation group morphism:
  From: Cyclic group of order 12 as a permutation group
  To:   Cyclic group of order 20 as a permutation group
  Defn: [(1,2,3,4,5,6,7,8,9,10,11,12)] ->
        [(1,6,11,16)(2,7,12,17)(3,8,13,18)(4,9,14,19)(5,10,15,20)]
```

```
a = C12("(1,6,11,4,9,2,7,12,5,10,3,8)")
phi(a)
```

```
(1,6,11,16)(2,7,12,17)(3,8,13,18)(4,9,14,19)(5,10,15,20)
```

```
b = C12("(1,3,5,7,9,11)(2,4,6,8,10,12)")
phi(b)
```

```
(1,11)(2,12)(3,13)(4,14)(5,15)(6,16)(7,17)(8,18)(9,19)(10,20)
```

```
c = C12("(1,9,5)(2,10,6)(3,11,7)(4,12,8)")
phi(c)
```

```
()
```

Note that the element c must therefore be in the kernel of phi.

   We can then compute the subgroup of the domain that is the kernel, in this case a cyclic group of order 3 inside the cyclic group of order 12. We can compute the image of *any* subgroup, but here we will build the whole homomorphic image by supplying the whole domain to the .image() method. Here the image is a cyclic subgroup of order 4 inside the cyclic group of order 20. Then we can verify the First Isomorphism Theorem.

```
K = phi.kernel(); K
```

```
Subgroup of (Cyclic group of order 12 as a permutation group)
generated by [(1,5,9)(2,6,10)(3,7,11)(4,8,12)]
```

```
Im = phi.image(C12); Im
```

```
Subgroup of (Cyclic group of order 20 as a permutation group)
generated by [(1,6,11,16)(2,7,12,17)(3,8,13,18)(4,9,14,19)(5,10,15,20)]
```

```
Im.is_isomorphic(C12.quotient(K))
```

```
True
```

Here is a slightly more complicated example. The dihedral group D_{20} is the symmetry group of a 20-gon. Inside this group is a subgroup that is isomorphic to the symmetry group of a 5-gon (pentagon). Is this a surprise, or is this obvious? Here is a way to make precise the statement "D_{20} contains a copy of D_5."

   We build the domain and find its generators, so we know how many images to supply in the definition of the homomorphism. Then we construct the codomain, from which we will construct images. Our choice here is to send a reflection to a reflection, and a rotation to a rotation. But the rotations will both have order 5, and both are a rotation by 72 degrees.

```
G = DihedralGroup(5)
H = DihedralGroup(20)
G.gens()
```

```
[(1,2,3,4,5), (1,5)(2,4)]
```

```
H.gens()
```

```
[(1,2,3,4,5,6,7,8,9,10,11,12,13,14,15,16,17,18,19,20),
 (1,20)(2,19)(3,18)(4,17)(5,16)(6,15)(7,14)(8,13)(9,12)(10,11)]
```

```
x = H.gen(0)^4
y = H.gen(1)
rho = PermutationGroupMorphism(G, H, [x, y])
rho.kernel()
```

```
Subgroup of (Dihedral group of order 10 as a permutation group)
generated by [()]
```

Since the kernel is trivial, `rho` is a one-to-one function (see Exercise 11.3.18). But more importantly, by the First Isomorphishm Theorem, `G` is isomorphic to the image of the homomorphism. We compute the image and check the claim.

```
Im = rho.image(G); Im
```

```
Subgroup of (Dihedral group of order 40 as a permutation group)
generated by
[(1,5,9,13,17)(2,6,10,14,18)(3,7,11,15,19)(4,8,12,16,20),
(1,20)(2,19)(3,18)(4,17)(5,16)(6,15)(7,14)(8,13)(9,12)(10,11)]
```

```
Im.is_subgroup(H)
```

```
True
```

```
Im.is_isomorphic(G)
```

```
True
```

Just providing a list of images for the generators of the domain is no guarantee that the function will extend to a homomorphism. For starters, the order of each image must divide the order of the corresponding preimage. (Can you prove this?) And similarly, if the domain is abelian, then the image must also be abelian, so in this case the list of images should not generate a non-abelian subgroup. Here is an example. There are no homomorphisms from a cyclic group of order 7 to a cyclic group of order 4 (other than the trivial function that takes every element to the identity). To see this, consider the possible orders of the kernel, and of the two possibilities, see that one is impossible and the other arises with the trivial homomorphism. Unfortunately, Sage acts as if nothing is wrong in creating a homomorphism between these groups, but what Sage builds is useless and raises errors when you try to use it.

```
G = CyclicPermutationGroup(7)
H = CyclicPermutationGroup(4)
tau = PermutationGroupMorphism_im_gens(G, H, H.gens())
tau
```

```
Permutation group morphism:
  From: Cyclic group of order 7 as a permutation group
  To:   Cyclic group of order 4 as a permutation group
  Defn: [(1,2,3,4,5,6,7)] -> [(1,2,3,4)]
```

```
tau.kernel()
```

```
Traceback (most recent call last):
...
RuntimeError: Gap produced error output
...
```

Rather than creating homomorphisms ourselves, in certain situations Sage knows of the existence of natural homomorphisms and will create them for you. One such case is a direct product construction. Given a group G, the method `.direct_product(H)` will create the direct product $G \times H$. (This is not the same command as the function `direct_product_permgroups()` from before.) Not only does this command create the direct product, but it also builds *four* homomorphisms, one with domain G, one with domain H and two with domain $G \times H$. So the output consists of five objects, the first being the actual group, and the remainder are homomorphisms. We will demonstrate the call here, and leave a more thorough investigation for the exercises.

```
G = CyclicPermutationGroup(3)
H = DihedralGroup(4)
results = G.direct_product(H)
results[0]
```

```
Permutation Group with generators [(4,5,6,7), (4,7)(5,6), (1,2,3)]
```

```
results[1]
```

```
Permutation group morphism:
  From: Cyclic group of order 3 as a permutation group
  To:   Permutation Group with generators
        [(4,5,6,7), (4,7)(5,6), (1,2,3)]
  Defn: Embedding( Group( [ (1,2,3), (4,5,6,7), (4,7)(5,6) ] ), 1 )
```

```
results[2]
```

```
Permutation group morphism:
  From: Dihedral group of order 8 as a permutation group
  To:   Permutation Group with generators
        [(4,5,6,7), (4,7)(5,6), (1,2,3)]
  Defn: Embedding( Group( [ (1,2,3), (4,5,6,7), (4,7)(5,6) ] ), 2 )
```

```
results[3]
```

```
Permutation group morphism:
  From: Permutation Group with generators
        [(4,5,6,7), (4,7)(5,6), (1,2,3)]
  To:   Cyclic group of order 3 as a permutation group
  Defn: Projection( Group( [ (1,2,3), (4,5,6,7), (4,7)(5,6) ] ), 1 )
```

```
results[4]
```

```
Permutation group morphism:
  From: Permutation Group with generators
        [(4,5,6,7), (4,7)(5,6), (1,2,3)]
  To:   Dihedral group of order 8 as a permutation group
  Defn: Projection( Group( [ (1,2,3), (4,5,6,7), (4,7)(5,6) ] ), 2 )
```

11.6   Sage Exercises

1.   An automorphism is an isomorphism between a group and itself. The identity function $(x \mapsto x)$ is always an isomorphism, which we consider trivial. Use Sage to construct a nontrivial automorphism of the cyclic group of order 12. Check that the mapping is both onto and one-to-one by computing the image and kernel and performing the proper tests on these subgroups. Now construct all of the possible automorphisms of the cyclic group of order 12 without any duplicates.

2.   The four homomorphisms created by the direct product construction are each an example of a more general construction of homomorphisms involving groups G, H and $G \times H$. By using the same groups as in the example in the previous subsection, see if you can discover and describe these constructions with exact definitions of the four homomorphisms in general.

Your tools for investigating a Sage group homomorphism are limited, you might take each generator of the domain and see what its image is. Here is an example of the type of computation you might do repeatedly. We'll investigate the second homomorphism. The domain is the dihedral group, and we will compute the image of the first generator.

```
G = CyclicPermutationGroup(3)
H = DihedralGroup(4)
results = G.direct_product(H)
phi = results[2]
H.gens()
```

```
[(1,2,3,4), (1,4)(2,3)]
```

```
a = H.gen(0); a
```

```
(1,2,3,4)
```

```
phi(a)
```

```
(4,5,6,7)
```

3.   Consider two permutation groups. The first is the subgroup of S_7 generated by $(1, 2, 3)$ and $(4, 5, 6, 7)$. The second is a subgroup of S_{12} generated by $(1, 2, 3)(4, 5, 6)(7, 8, 9)(10, 11, 12)$ and $(1, 10, 7, 4)(2, 11, 8, 5)(3, 12, 9, 6)$. Build these two groups and use the proper Sage command to see that they are isomorphic. Then construct a homomorphism between these two groups that is an isomorphism and include enough details to verify that the mapping is really an isomorphism.

4.   The second paragraph of this chapter informally describes a homomorphism from S_n to $\mathbb{Z}_2$, where the even permutations all map to one of the elements and the odd permutations all map to the other element. Replace S_n by S_6 and replace $\mathbb{Z}_2$ by the permutation version of the cyclic subgroup of order 2, and construct a nontrivial homomorphism between these two groups. Evaluate your homomorphism with enough even and odd permutations to be convinced that it is correct. Then construct the kernel and verify that it is the group you expect.

Hints: First, decide which elements of the group of order 2 will be associated with even permutations and which will be associated with odd permutations. Then examine the generators of S_6 to help decide just how to build the homomorphism.

5.   The dihedral group D_{20} has several normal subgroups, as seen below. Each of these is the kernel of a homomorphism with D_{20} as the domain. For each normal subgroup of D_{20} construct a homomorphism from D_{20} to D_{20} that has the normal subgroup as the kernel.

Include in your work verifications that you are creating the desired kernels. There is a pattern to many of these, but the three of order 20 will be a challenge.

```
G = DihedralGroup(20)
[H.order() for H in G.normal_subgroups()]
```

```
[1, 2, 4, 5, 10, 20, 20, 20, 40]
```

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<table><tr><td>12</td></tr></table>
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Matrix Groups and Symmetry

When Felix Klein (1849–1925) accepted a chair at the University of Erlangen, he outlined in his inaugural address a program to classify different geometries. Central to Klein's program was the theory of groups: he considered geometry to be the study of properties that are left invariant under transformation groups. Groups, especially matrix groups, have now become important in the study of symmetry and have found applications in such disciplines as chemistry and physics. In the first part of this chapter, we will examine some of the classical matrix groups, such as the general linear group, the special linear group, and the orthogonal group. We will then use these matrix groups to investigate some of the ideas behind geometric symmetry.

12.1 Matrix Groups

Some Facts from Linear Algebra

Before we study matrix groups, we must recall some basic facts from linear algebra. One of the most fundamental ideas of linear algebra is that of a linear transformation. A **_linear transformation_** or **_linear map_** $T : \mathbb{R}^n \to \mathbb{R}^m$ is a map that preserves vector addition and scalar multiplication; that is, for vectors $\mathbf{x}$ and $\mathbf{y}$ in $\mathbb{R}^n$ and a scalar $\alpha \in \mathbb{R}$,

$$T(\mathbf{x} + \mathbf{y}) = T(\mathbf{x}) + T(\mathbf{y})$$
$$T(\alpha\mathbf{y}) = \alpha T(\mathbf{y}).$$

An $m \times n$ matrix with entries in $\mathbb{R}$ represents a linear transformation from $\mathbb{R}^n$ to $\mathbb{R}^m$. If we write vectors $\mathbf{x} = (x_1, \ldots, x_n)^{\mathrm{t}}$ and $\mathbf{y} = (y_1, \ldots, y_n)^{\mathrm{t}}$ in $\mathbb{R}^n$ as column matrices, then an $m \times n$ matrix

$$A = \begin{pmatrix} a_{11} & a_{12} & \cdots & a_{1n} \\ a_{21} & a_{22} & \cdots & a_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ a_{m1} & a_{m2} & \cdots & a_{mn} \end{pmatrix}$$

maps the vectors to $\mathbb{R}^m$ linearly by matrix multiplication. Observe that if α is a real number,

$$A(\mathbf{x} + \mathbf{y}) = A\mathbf{x} + A\mathbf{y} \qquad \text{and} \qquad \alpha A\mathbf{x} = A(\alpha\mathbf{x}),$$

where

$$\mathbf{x} = \begin{pmatrix} x_1 \\ x_2 \\ \vdots \\ x_n \end{pmatrix}.$$

<div align="center">194</div>

We will often abbreviate the matrix A by writing (a_{ij}).

Conversely, if $T : \mathbb{R}^n \to \mathbb{R}^m$ is a linear map, we can associate a matrix A with T by considering what T does to the vectors

$$
\begin{aligned}
\mathbf{e}_1 &= (1, 0, \ldots, 0)^{\mathrm{t}} \\
\mathbf{e}_2 &= (0, 1, \ldots, 0)^{\mathrm{t}} \\
&\vdots \\
\mathbf{e}_n &= (0, 0, \ldots, 1)^{\mathrm{t}}.
\end{aligned}
$$

We can write any vector $\mathbf{x} = (x_1, \ldots, x_n)^{\mathrm{t}}$ as

$$
x_1 \mathbf{e}_1 + x_2 \mathbf{e}_2 + \cdots + x_n \mathbf{e}_n.
$$

Consequently, if

$$
\begin{aligned}
T(\mathbf{e}_1) &= (a_{11}, a_{21}, \ldots, a_{m1})^{\mathrm{t}}, \\
T(\mathbf{e}_2) &= (a_{12}, a_{22}, \ldots, a_{m2})^{\mathrm{t}}, \\
&\vdots \\
T(\mathbf{e}_n) &= (a_{1n}, a_{2n}, \ldots, a_{mn})^{\mathrm{t}},
\end{aligned}
$$

then

$$
\begin{aligned}
T(\mathbf{x}) &= T(x_1 \mathbf{e}_1 + x_2 \mathbf{e}_2 + \cdots + x_n \mathbf{e}_n) \\
&= x_1 T(\mathbf{e}_1) + x_2 T(\mathbf{e}_2) + \cdots + x_n T(\mathbf{e}_n) \\
&= \left(\sum_{k=1}^{n} a_{1k} x_k, \ldots, \sum_{k=1}^{n} a_{mk} x_k \right)^{\mathrm{t}} \\
&= A\mathbf{x}.
\end{aligned}
$$

Example 12.1. If we let $T : \mathbb{R}^2 \to \mathbb{R}^2$ be the map given by

$$
T(x_1, x_2) = (2x_1 + 5x_2, -4x_1 + 3x_2),
$$

the axioms that T must satisfy to be a linear transformation are easily verified. The column vectors $T\mathbf{e}_1 = (2, -4)^{\mathrm{t}}$ and $T\mathbf{e}_2 = (5, 3)^{\mathrm{t}}$ tell us that T is given by the matrix

$$
A = \begin{pmatrix} 2 & 5 \\ -4 & 3 \end{pmatrix}.
$$

Since we are interested in groups of matrices, we need to know which matrices have multiplicative inverses. Recall that an $n \times n$ matrix A is **invertible** exactly when there exists another matrix A^{-1} such that $AA^{-1} = A^{-1}A = I$, where

$$
I = \begin{pmatrix} 1 & 0 & \cdots & 0 \\ 0 & 1 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & 1 \end{pmatrix}
$$

is the $n \times n$ identity matrix. From linear algebra we know that A is invertible if and only if the determinant of A is nonzero. Sometimes an invertible matrix is said to be **nonsingular**.

Example 12.2. If A is the matrix

$$\begin{pmatrix} 2 & 1 \\ 5 & 3 \end{pmatrix},$$

then the inverse of A is

$$A^{-1} = \begin{pmatrix} 3 & -1 \\ -5 & 2 \end{pmatrix}.$$

We are guaranteed that A^{-1} exists, since $\det(A) = 2 \cdot 3 - 5 \cdot 1 = 1$ is nonzero.

Some other facts about determinants will also prove useful in the course of this chapter. Let A and B be $n \times n$ matrices. From linear algebra we have the following properties of determinants.

- The determinant is a homomorphism into the multiplicative group of real numbers; that is, $\det(AB) = (\det A)(\det B)$.

- If A is an invertible matrix, then $\det(A^{-1}) = 1/\det A$.

- If we define the transpose of a matrix $A = (a_{ij})$ to be $A^{\mathrm{t}} = (a_{ji})$, then $\det(A^{\mathrm{t}}) = \det A$.

- Let T be the linear transformation associated with an $n \times n$ matrix A. Then T multiplies volumes by a factor of $|\det A|$. In the case of $\mathbb{R}^2$, this means that T multiplies areas by $|\det A|$.

Linear maps, matrices, and determinants are covered in any elementary linear algebra text; however, if you have not had a course in linear algebra, it is a straightforward process to verify these properties directly for 2×2 matrices, the case with which we are most concerned.

The General and Special Linear Groups

The set of all $n \times n$ invertible matrices forms a group called the **general linear group**. We will denote this group by $GL_n(\mathbb{R})$. The general linear group has several important subgroups. The multiplicative properties of the determinant imply that the set of matrices with determinant one is a subgroup of the general linear group. Stated another way, suppose that $\det(A) = 1$ and $\det(B) = 1$. Then $\det(AB) = \det(A)\det(B) = 1$ and $\det(A^{-1}) = 1/\det A = 1$. This subgroup is called the **special linear group** and is denoted by $SL_n(\mathbb{R})$.

Example 12.3. Given a 2×2 matrix

$$A = \begin{pmatrix} a & b \\ c & d \end{pmatrix},$$

the determinant of A is $ad - bc$. The group $GL_2(\mathbb{R})$ consists of those matrices in which $ad - bc \neq 0$. The inverse of A is

$$A^{-1} = \frac{1}{ad - bc} \begin{pmatrix} d & -b \\ -c & a \end{pmatrix}.$$

If A is in $SL_2(\mathbb{R})$, then

$$A^{-1} = \begin{pmatrix} d & -b \\ -c & a \end{pmatrix}.$$

Geometrically, $SL_2(\mathbb{R})$ is the group that preserves the areas of parallelograms. Let

$$A = \begin{pmatrix} 1 & 1 \\ 0 & 1 \end{pmatrix}$$

be in $SL_2(\mathbb{R})$. In Figure 12.4, the unit square corresponding to the vectors $\mathbf{x} = (1, 0)^t$ and $\mathbf{y} = (0, 1)^t$ is taken by A to the parallelogram with sides $(1, 0)^t$ and $(1, 1)^t$; that is, $A\mathbf{x} = (1, 0)^t$ and $A\mathbf{y} = (1, 1)^t$. Notice that these two parallelograms have the same area.
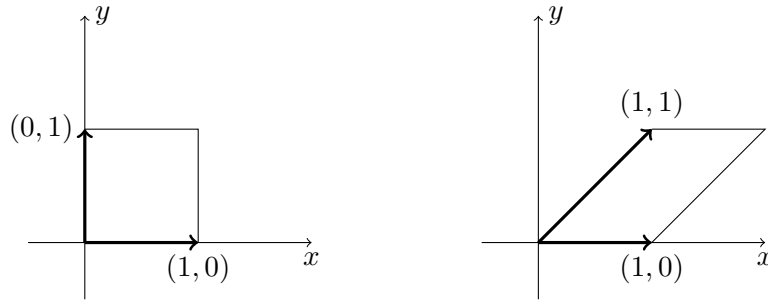

Figure 12.4: $SL_2(\mathbb{R})$ acting on the unit square

The Orthogonal Group $O(n)$

Another subgroup of $GL_n(\mathbb{R})$ is the orthogonal group. A matrix A is ***orthogonal*** if $A^{-1} = A^t$. The ***orthogonal group*** consists of the set of all orthogonal matrices. We write $O(n)$ for the $n \times n$ orthogonal group. We leave as an exercise the proof that $O(n)$ is a subgroup of $GL_n(\mathbb{R})$.

Example 12.5. The following matrices are orthogonal:

$$\begin{pmatrix} 3/5 & -4/5 \\ 4/5 & 3/5 \end{pmatrix}, \quad \begin{pmatrix} 1/2 & -\sqrt{3}/2 \\ \sqrt{3}/2 & 1/2 \end{pmatrix}, \quad \begin{pmatrix} -1/\sqrt{2} & 0 & 1/\sqrt{2} \\ 1/\sqrt{6} & -2/\sqrt{6} & 1/\sqrt{6} \\ 1/\sqrt{3} & 1/\sqrt{3} & 1/\sqrt{3} \end{pmatrix}.$$

There is a more geometric way of viewing the group $O(n)$. The orthogonal matrices are exactly those matrices that preserve the length of vectors. We can define the length of a vector using the ***Euclidean inner product***, or ***dot product***, of two vectors. The Euclidean inner product of two vectors $\mathbf{x} = (x_1, \ldots, x_n)^t$ and $\mathbf{y} = (y_1, \ldots, y_n)^t$ is

$$\langle \mathbf{x}, \mathbf{y} \rangle = \mathbf{x}^t \mathbf{y} = (x_1, x_2, \ldots, x_n) \begin{pmatrix} y_1 \\ y_2 \\ \vdots \\ y_n \end{pmatrix} = x_1 y_1 + \cdots + x_n y_n.$$

We define the length of a vector $\mathbf{x} = (x_1, \ldots, x_n)^t$ to be

$$\|\mathbf{x}\| = \sqrt{\langle \mathbf{x}, \mathbf{x} \rangle} = \sqrt{x_1^2 + \cdots + x_n^2}.$$

Associated with the notion of the length of a vector is the idea of the distance between two vectors. We define the ***distance*** between two vectors $\mathbf{x}$ and $\mathbf{y}$ to be $\|\mathbf{x} - \mathbf{y}\|$. We leave as an exercise the proof of the following proposition about the properties of Euclidean inner products.

Proposition 12.6. *Let* **x**, **y**, *and* **w** *be vectors in* $\mathbb{R}^n$ *and* $\alpha \in \mathbb{R}$. *Then*

1. $\langle \mathbf{x}, \mathbf{y} \rangle = \langle \mathbf{y}, \mathbf{x} \rangle$.

2. $\langle \mathbf{x}, \mathbf{y} + \mathbf{w} \rangle = \langle \mathbf{x}, \mathbf{y} \rangle + \langle \mathbf{x}, \mathbf{w} \rangle$.

3. $\langle \alpha \mathbf{x}, \mathbf{y} \rangle = \langle \mathbf{x}, \alpha \mathbf{y} \rangle = \alpha \langle \mathbf{x}, \mathbf{y} \rangle$.

4. $\langle \mathbf{x}, \mathbf{x} \rangle \geq 0$ *with equality exactly when* $\mathbf{x} = 0$.

5. *If* $\langle \mathbf{x}, \mathbf{y} \rangle = 0$ *for all* **x** *in* $\mathbb{R}^n$, *then* $\mathbf{y} = 0$.

Example 12.7. The vector $\mathbf{x} = (3, 4)^{\text{t}}$ has length $\sqrt{3^2 + 4^2} = 5$. We can also see that the orthogonal matrix

$$A = \begin{pmatrix} 3/5 & -4/5 \\ 4/5 & 3/5 \end{pmatrix}$$

preserves the length of this vector. The vector $A\mathbf{x} = (-7/5, 24/5)^{\text{t}}$ also has length 5.

Since $\det(AA^{\text{t}}) = \det(I) = 1$ and $\det(A) = \det(A^{\text{t}})$, the determinant of any orthogonal matrix is either 1 or -1. Consider the column vectors

$$\mathbf{a}_j = \begin{pmatrix} a_{1j} \\ a_{2j} \\ \vdots \\ a_{nj} \end{pmatrix}$$

of the orthogonal matrix $A = (a_{ij})$. Since $AA^{\text{t}} = I$, $\langle \mathbf{a}_r, \mathbf{a}_s \rangle = \delta_{rs}$, where

$$\delta_{rs} = \begin{cases} 1 & r = s \\ 0 & r \neq s \end{cases}$$

is the Kronecker delta. Accordingly, column vectors of an orthogonal matrix all have length 1; and the Euclidean inner product of distinct column vectors is zero. Any set of vectors satisfying these properties is called an ***orthonormal set***. Conversely, given an $n \times n$ matrix A whose columns form an orthonormal set, it follows that $A^{-1} = A^{\text{t}}$.

We say that a matrix A is ***distance-preserving***, ***length-preserving***, or ***inner product-preserving*** when $\|T\mathbf{x} - T\mathbf{y}\| = \|\mathbf{x} - \mathbf{y}\|$, $\|T\mathbf{x}\| = \|\mathbf{x}\|$, or $\langle T\mathbf{x}, T\mathbf{y} \rangle = \langle \mathbf{x}, \mathbf{y} \rangle$, respectively. The following theorem, which characterizes the orthogonal group, says that these notions are the same.

Theorem 12.8. *Let A be an $n \times n$ matrix. The following statements are equivalent.*

1. *The columns of the matrix A form an orthonormal set.*

2. $A^{-1} = A^{\text{t}}$.

3. *For vectors* **x** *and* **y**, $\langle A\mathbf{x}, A\mathbf{y} \rangle = \langle \mathbf{x}, \mathbf{y} \rangle$.

4. *For vectors* **x** *and* **y**, $\|A\mathbf{x} - A\mathbf{y}\| = \|\mathbf{x} - \mathbf{y}\|$.

5. *For any vector* **x**, $\|A\mathbf{x}\| = \|\mathbf{x}\|$.

PROOF.   We have already shown (1) and (2) to be equivalent.
   $(2) \Rightarrow (3)$.

$$\langle A\mathbf{x}, A\mathbf{y} \rangle = (A\mathbf{x})^{\mathrm{t}} A\mathbf{y}$$
$$= \mathbf{x}^{\mathrm{t}} A^{\mathrm{t}} A\mathbf{y}$$
$$= \mathbf{x}^{\mathrm{t}} \mathbf{y}$$
$$= \langle \mathbf{x}, \mathbf{y} \rangle.$$

$(3) \Rightarrow (2)$. Since

$$\langle \mathbf{x}, \mathbf{x} \rangle = \langle A\mathbf{x}, A\mathbf{x} \rangle$$
$$= \mathbf{x}^{\mathrm{t}} A^{\mathrm{t}} A\mathbf{x}$$
$$= \langle \mathbf{x}, A^{\mathrm{t}} A\mathbf{x} \rangle,$$

we know that $\langle \mathbf{x}, (A^{\mathrm{t}} A - I)\mathbf{x} \rangle = 0$ for all $\mathbf{x}$. Therefore, $A^{\mathrm{t}} A - I = 0$ or $A^{-1} = A^{\mathrm{t}}$.
   $(3) \Rightarrow (4)$. If A is inner product-preserving, then A is distance-preserving, since

$$\| A\mathbf{x} - A\mathbf{y} \|^2 = \| A(\mathbf{x} - \mathbf{y}) \|^2$$
$$= \langle A(\mathbf{x} - \mathbf{y}), A(\mathbf{x} - \mathbf{y}) \rangle$$
$$= \langle \mathbf{x} - \mathbf{y}, \mathbf{x} - \mathbf{y} \rangle$$
$$= \| \mathbf{x} - \mathbf{y} \|^2.$$

$(4) \Rightarrow (5)$. If A is distance-preserving, then A is length-preserving. Letting $\mathbf{y} = 0$, we have

$$\| A\mathbf{x} \| = \| A\mathbf{x} - A\mathbf{y} \| = \| \mathbf{x} - \mathbf{y} \| = \| \mathbf{x} \|.$$

$(5) \Rightarrow (3)$. We use the following identity to show that length-preserving implies inner product-preserving:

$$\langle \mathbf{x}, \mathbf{y} \rangle = \frac{1}{2} \left[\| \mathbf{x} + \mathbf{y} \|^2 - \| \mathbf{x} \|^2 - \| \mathbf{y} \|^2 \right].$$

Observe that

$$\langle A\mathbf{x}, A\mathbf{y} \rangle = \frac{1}{2} \left[\| A\mathbf{x} + A\mathbf{y} \|^2 - \| A\mathbf{x} \|^2 - \| A\mathbf{y} \|^2 \right]$$
$$= \frac{1}{2} \left[\| A(\mathbf{x} + \mathbf{y}) \|^2 - \| A\mathbf{x} \|^2 - \| A\mathbf{y} \|^2 \right]$$
$$= \frac{1}{2} \left[\| \mathbf{x} + \mathbf{y} \|^2 - \| \mathbf{x} \|^2 - \| \mathbf{y} \|^2 \right]$$
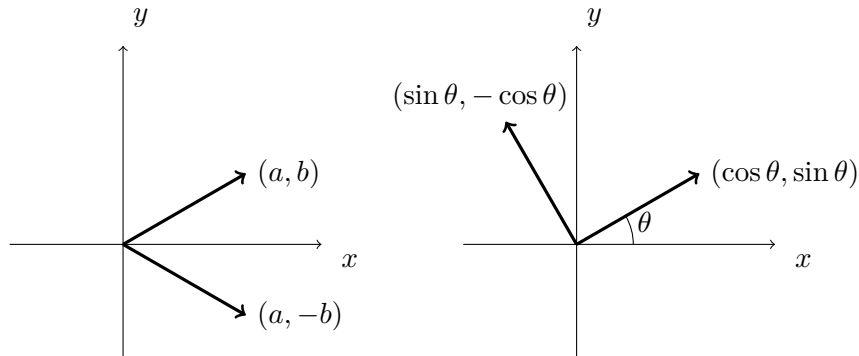$$= \langle \mathbf{x}, \mathbf{y} \rangle. \qquad \square$$

Figure 12.9: $O(2)$ acting on $\mathbb{R}^2$

Example 12.10. Let us examine the orthogonal group on $\mathbb{R}^2$ a bit more closely. An element $T \in O(2)$ is determined by its action on $\mathbf{e}_1 = (1,0)^{\mathrm{t}}$ and $\mathbf{e}_2 = (0,1)^{\mathrm{t}}$. If $T(\mathbf{e}_1) = (a,b)^{\mathrm{t}}$, then $a^2 + b^2 = 1$ and $T(\mathbf{e}_2) = (-b,a)^{\mathrm{t}}$. Hence, T can be represented by

$$A = \begin{pmatrix} a & -b \\ b & a \end{pmatrix} = \begin{pmatrix} \cos\theta & -\sin\theta \\ \sin\theta & \cos\theta \end{pmatrix},$$

where $0 \leq \theta < 2\pi$. A matrix T in $O(2)$ either reflects or rotates a vector in $\mathbb{R}^2$ (Figure 12.9). A reflection about the horizontal axis is given by the matrix

$$\begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix},$$

whereas a rotation by an angle θ in a counterclockwise direction must come from a matrix of the form

$$\begin{pmatrix} \cos\theta & \sin\theta \\ \sin\theta & -\cos\theta \end{pmatrix}.$$

A reflection about a line ℓ is simply a reflection about the horizontal axis followed by a rotation. If $\det A = -1$, then A gives a reflection.

Two of the other matrix or matrix-related groups that we will consider are the special orthogonal group and the group of Euclidean motions. The **special orthogonal group**, $SO(n)$, is just the intersection of $O(n)$ and $SL_n(\mathbb{R})$; that is, those elements in $O(n)$ with determinant one. The **Euclidean group**, $E(n)$, can be written as ordered pairs $(A, \mathbf{x})$, where A is in $O(n)$ and $\mathbf{x}$ is in $\mathbb{R}^n$. We define multiplication by

$$(A, \mathbf{x})(B, \mathbf{y}) = (AB, A\mathbf{y} + \mathbf{x}).$$

The identity of the group is $(I, \mathbf{0})$; the inverse of $(A, \mathbf{x})$ is $(A^{-1}, -A^{-1}\mathbf{x})$. In Exercise 12.3.6, you are asked to check that $E(n)$ is indeed a group under this operation.
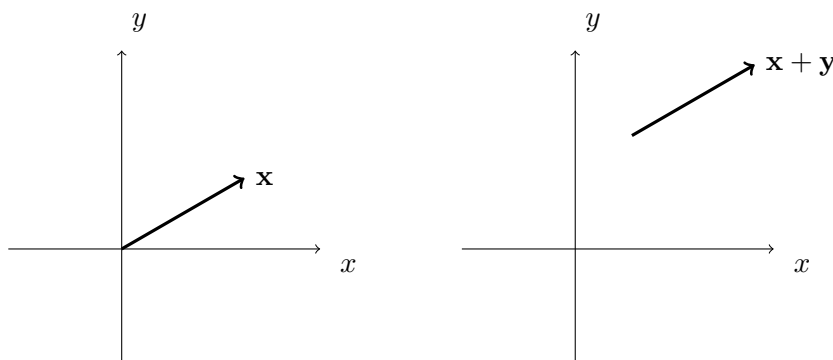

Figure 12.11: Translations in $\mathbb{R}^2$

12.2   Symmetry

An **isometry** or **rigid motion** in $\mathbb{R}^n$ is a distance-preserving function f from $\mathbb{R}^n$ to $\mathbb{R}^n$. This means that f must satisfy

$$\|f(\mathbf{x}) - f(\mathbf{y})\| = \|\mathbf{x} - \mathbf{y}\|$$

for all $\mathbf{x}, \mathbf{y} \in \mathbb{R}^n$. It is not difficult to show that f must be a one-to-one map. By Theorem 12.8, any element in $O(n)$ is an isometry on $\mathbb{R}^n$; however, $O(n)$ does not include all

possible isometries on $\mathbb{R}^n$. Translation by a vector $\mathbf{x}$, $T_{\mathbf{y}}(\mathbf{x}) = \mathbf{x} + \mathbf{y}$ is also an isometry (Figure 12.11); however, T cannot be in $O(n)$ since it is not a linear map.

We are mostly interested in isometries in $\mathbb{R}^2$. In fact, the only isometries in $\mathbb{R}^2$ are rotations and reflections about the origin, translations, and combinations of the two. For example, a **_glide reflection_** is a translation followed by a reflection (Figure 12.12). In $\mathbb{R}^n$ all isometries are given in the same manner. The proof is very easy to generalize.
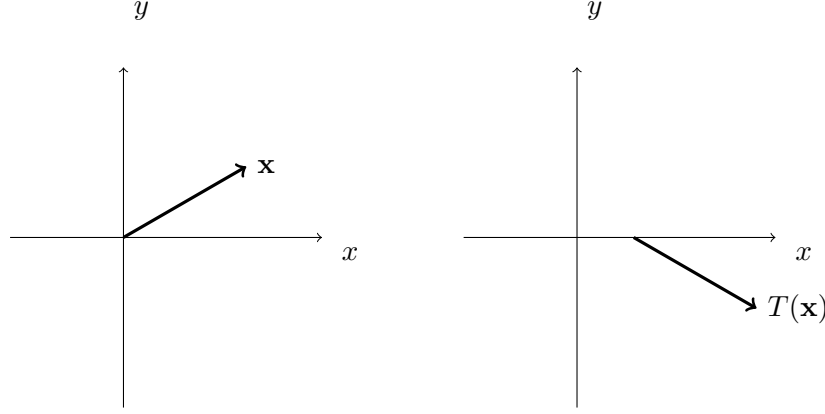
Figure 12.12: Glide reflections

Lemma 12.13. *An isometry f that fixes the origin in $\mathbb{R}^2$ is a linear transformation. In particular, f is given by an element in $O(2)$.*

PROOF. Let f be an isometry in $\mathbb{R}^2$ fixing the origin. We will first show that f preserves inner products. Since $f(0) = 0$, $\|f(\mathbf{x})\| = \|\mathbf{x}\|$; therefore,

$$
\begin{aligned}
\|\mathbf{x}\|^2 - 2\langle f(\mathbf{x}), f(\mathbf{y})\rangle + \|\mathbf{y}\|^2 &= \|f(\mathbf{x})\|^2 - 2\langle f(\mathbf{x}), f(\mathbf{y})\rangle + \|f(\mathbf{y})\|^2 \\
&= \langle f(\mathbf{x}) - f(\mathbf{y}), f(\mathbf{x}) - f(\mathbf{y})\rangle \\
&= \|f(\mathbf{x}) - f(\mathbf{y})\|^2 \\
&= \|\mathbf{x} - \mathbf{y}\|^2 \\
&= \langle \mathbf{x} - \mathbf{y}, \mathbf{x} - \mathbf{y}\rangle \\
&= \|\mathbf{x}\|^2 - 2\langle \mathbf{x}, \mathbf{y}\rangle + \|\mathbf{y}\|^2.
\end{aligned}
$$

Consequently,

$$\langle f(\mathbf{x}), f(\mathbf{y})\rangle = \langle \mathbf{x}, \mathbf{y}\rangle.$$

Now let $\mathbf{e}_1$ and $\mathbf{e}_2$ be $(1,0)^{\mathrm{t}}$ and $(0,1)^{\mathrm{t}}$, respectively. If

$$\mathbf{x} = (x_1, x_2) = x_1\mathbf{e}_1 + x_2\mathbf{e}_2,$$

then

$$f(\mathbf{x}) = \langle f(\mathbf{x}), f(\mathbf{e}_1)\rangle f(\mathbf{e}_1) + \langle f(\mathbf{x}), f(\mathbf{e}_2)\rangle f(\mathbf{e}_2) = x_1 f(\mathbf{e}_1) + x_2 f(\mathbf{e}_2).$$

The linearity of f easily follows. $\qquad\square$

For any arbitrary isometry, f, $T_{\mathbf{x}}f$ will fix the origin for some vector $\mathbf{x}$ in $\mathbb{R}^2$; hence, $T_{\mathbf{x}}f(\mathbf{y}) = A\mathbf{y}$ for some matrix $A \in O(2)$. Consequently, $f(\mathbf{y}) = A\mathbf{y} + \mathbf{x}$. Given the isometries

$$
\begin{aligned}
f(\mathbf{y}) &= A\mathbf{y} + \mathbf{x}_1 \\
g(\mathbf{y}) &= B\mathbf{y} + \mathbf{x}_2,
\end{aligned}
$$

their composition is

$$f(g(\mathbf{y})) = f(B\mathbf{y} + \mathbf{x}_2) = AB\mathbf{y} + A\mathbf{x}_2 + \mathbf{x}_1.$$

This last computation allows us to identify the group of isometries on $\mathbb{R}^2$ with $E(2)$.

Theorem 12.14. *The group of isometries on $\mathbb{R}^2$ is the Euclidean group, $E(2)$.*

A **symmetry group** in $\mathbb{R}^n$ is a subgroup of the group of isometries on $\mathbb{R}^n$ that fixes a set of points $X \subset \mathbb{R}^2$. It is important to realize that the symmetry group of X depends *both* on $\mathbb{R}^n$ and on X. For example, the symmetry group of the origin in $\mathbb{R}^1$ is $\mathbb{Z}_2$, but the symmetry group of the origin in $\mathbb{R}^2$ is $O(2)$.

Theorem 12.15. *The only finite symmetry groups in $\mathbb{R}^2$ are $\mathbb{Z}_n$ and D_n.*

PROOF.   Let $G = \{f_1, f_2, \ldots, f_n\}$ be a finite symmetry group that fixes a set of points in $X \subset \mathbb{R}^2$. Choose a point $\mathbf{x} \in X$. This point may not be a fixed point—it could be moved by G to another point in X. Define a set $S = \{\mathbf{y}_1, \mathbf{y}_2, \ldots \mathbf{y}_n\}$, where $\mathbf{y}_i = f_i(\mathbf{x})$. Now, let

$$\mathbf{z} = \frac{1}{n} \sum_{i=1}^{n} \mathbf{x}_i.$$

While the point $\mathbf{z}$ is not necessarily in the set X, it is fixed by every element in the symetry group. Without loss of generality, we may now assume that $\mathbf{z}$ is the origin.

Any finite symmetry group G in $\mathbb{R}^2$ that fixes the origin must be a finite subgroup of $O(2)$, since translations and glide reflections have infinite order. By Example 12.10, elements in $O(2)$ are either rotations of the form

$$R_\theta = \begin{pmatrix} \cos\theta & -\sin\theta \\ \sin\theta & \cos\theta \end{pmatrix}$$

or reflections of the form

$$T_\phi = \begin{pmatrix} \cos\phi & -\sin\phi \\ \sin\phi & \cos\phi \end{pmatrix} \begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix} = \begin{pmatrix} \cos\phi & \sin\phi \\ \sin\phi & -\cos\phi \end{pmatrix}.$$

Notice that $\det(R_\theta) = 1$, $\det(T_\phi) = -1$, and $T_\phi^2 = I$. We can divide the proof up into two cases. In the first case, all of the elements in G have determinant one. In the second case, there exists at least one element in G with determinant -1.

Case 1. The determinant of every element in G is one. In this case every element in G must be a rotation. Since G is finite, there is a smallest angle, say θ_0, such that the corresponding element R_{θ_0} is the smallest rotation in the positive direction. We claim that R_{θ_0} generates G. If not, then for some positive integer n there is an angle θ_1 between $n\theta_0$ and $(n+1)\theta_0$. If so, then $(n+1)\theta_0 - \theta_1$ corresponds to a rotation smaller than θ_0, which contradicts the minimality of θ_0.

Case 2. The group G contains a reflection T. The kernel of the homomorphism $\phi : G \to \{-1, 1\}$ given by $A \mapsto \det(A)$ consists of elements whose determinant is 1. Therefore, $|G/\ker\phi| = 2$. We know that the kernel is cyclic by the first case and is a subgroup of G of, say, order n. Hence, $|G| = 2n$. The elements of G are

$$R_\theta, \ldots, R_\theta^{n-1}, TR_\theta, \ldots, TR_\theta^{n-1}.$$

These elements satisfy the relation

$$TR_\theta T = R_\theta^{-1}.$$

Consequently, G must be isomorphic to D_n in this case. $\square$

The Wallpaper Groups

Suppose that we wish to study wallpaper patterns in the plane or crystals in three dimensions. Wallpaper patterns are simply repeating patterns in the plane (Figure 12.16). The analogs of wallpaper patterns in $\mathbb{R}^3$ are crystals, which we can think of as repeating patterns of molecules in three dimensions (Figure 12.17). The mathematical equivalent of a wallpaper or crystal pattern is called a lattice.
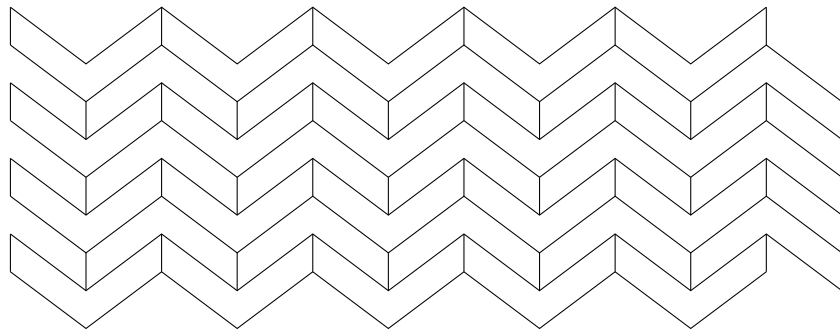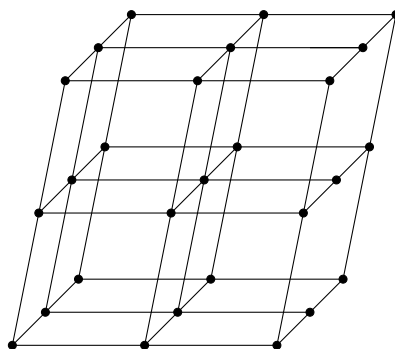

Figure 12.16: A wallpaper pattern in $\mathbb{R}^2$

Figure 12.17: A crystal structure in $\mathbb{R}^3$

Let us examine wallpaper patterns in the plane a little more closely. Suppose that $\mathbf{x}$ and $\mathbf{y}$ are linearly independent vectors in $\mathbb{R}^2$; that is, one vector cannot be a scalar multiple of the other. A **lattice** of $\mathbf{x}$ and $\mathbf{y}$ is the set of all linear combinations $m\mathbf{x} + n\mathbf{y}$, where m and n are integers. The vectors $\mathbf{x}$ and $\mathbf{y}$ are said to be a **basis** for the lattice.

Notice that a lattice can have several bases. For example, the vectors $(1,1)^{\mathrm{t}}$ and $(2,0)^{\mathrm{t}}$ have the same lattice as the vectors $(-1,1)^{\mathrm{t}}$ and $(-1,-1)^{\mathrm{t}}$ (Figure 12.18). However, any lattice is completely determined by a basis. Given two bases for the same lattice, say $\{\mathbf{x}_1, \mathbf{x}_2\}$ and $\{\mathbf{y}_1, \mathbf{y}_2\}$, we can write

$$\mathbf{y}_1 = \alpha_1 \mathbf{x}_1 + \alpha_2 \mathbf{x}_2$$
$$\mathbf{y}_2 = \beta_1 \mathbf{x}_1 + \beta_2 \mathbf{x}_2,$$

where α_1, α_2, β_1, and β_2 are integers. The matrix corresponding to this transformation is

$$U = \begin{pmatrix} \alpha_1 & \alpha_2 \\ \beta_1 & \beta_2 \end{pmatrix}.$$

If we wish to give $\mathbf{x}_1$ and $\mathbf{x}_2$ in terms of $\mathbf{y}_1$ and $\mathbf{y}_2$, we need only calculate U^{-1}; that is,

$$U^{-1} \begin{pmatrix} \mathbf{y}_1 \\ \mathbf{y}_2 \end{pmatrix} = \begin{pmatrix} \mathbf{x}_1 \\ \mathbf{x}_2 \end{pmatrix}.$$

Since U has integer entries, U^{-1} must also have integer entries; hence the determinants of both U and U^{-1} must be integers. Because $UU^{-1} = I$,

$$\det(UU^{-1}) = \det(U)\det(U^{-1}) = 1;$$

consequently, $\det(U) = \pm 1$. A matrix with determinant ± 1 and integer entries is called **unimodular**. For example, the matrix

$$\begin{pmatrix} 3 & 1 \\ 5 & 2 \end{pmatrix}$$

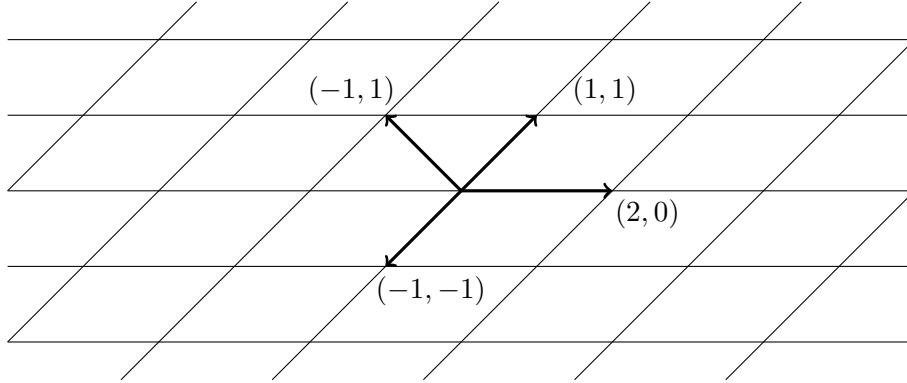is unimodular. It should be clear that there is a minimum length for vectors in a lattice.

Figure 12.18: A lattice in $\mathbb{R}^2$

We can classify lattices by studying their symmetry groups. The symmetry group of a lattice is the subgroup of $E(2)$ that maps the lattice to itself. We consider two lattices in $\mathbb{R}^2$ to be equivalent if they have the same symmetry group. Similarly, classification of crystals in $\mathbb{R}^3$ is accomplished by associating a symmetry group, called a **space group**, with each type of crystal. Two lattices are considered different if their space groups are not the same. The natural question that now arises is how many space groups exist.

A space group is composed of two parts: a **translation subgroup** and a **point**. The translation subgroup is an infinite abelian subgroup of the space group made up of the translational symmetries of the crystal; the point group is a finite group consisting of rotations and reflections of the crystal about a point. More specifically, a space group is a subgroup of $G \subset E(2)$ whose translations are a set of the form $\{(I, t) : t \in L\}$, where L is a lattice. Space groups are, of course, infinite. Using geometric arguments, we can prove the following theorem (see [5] or [6]).

Theorem 12.19. *Every translation group in $\mathbb{R}^2$ is isomorphic to $\mathbb{Z} \times \mathbb{Z}$.*

The point group of G is $G_0 = \{A : (A, b) \in G \text{ for some } b\}$. In particular, G_0 must be a subgroup of $O(2)$. Suppose that $\mathbf{x}$ is a vector in a lattice L with space group G, translation group H, and point group G_0. For any element $(A, \mathbf{y})$ in G,

$$\begin{aligned} (A, \mathbf{y})(I, \mathbf{x})(A, \mathbf{y})^{-1} &= (A, A\mathbf{x} + \mathbf{y})(A^{-1}, -A^{-1}\mathbf{y}) \\ &= (AA^{-1}, -AA^{-1}\mathbf{y} + A\mathbf{x} + \mathbf{y}) \\ &= (I, A\mathbf{x}); \end{aligned}$$

hence, $(I, A\mathbf{x})$ is in the translation group of G. More specifically, $A\mathbf{x}$ must be in the lattice L. It is important to note that G_0 is not usually a subgroup of the space group G; however, if T is the translation subgroup of G, then $G/T \cong G_0$. The proof of the following theorem can be found in [2], [5], or [6].

Theorem 12.20. *The point group in the wallpaper groups is isomorphic to $\mathbb{Z}_n$ or D_n, where $n = 1, 2, 3, 4, 6$.*

To answer the question of how the point groups and the translation groups can be combined, we must look at the different types of lattices. Lattices can be classified by the structure of a single lattice cell. The possible cell shapes are parallelogram, rectangular, square, rhombic, and hexagonal (Figure 12.21). The wallpaper groups can now be classified according to the types of reflections that occur in each group: these are ordinarily reflections, glide reflections, both, or none.
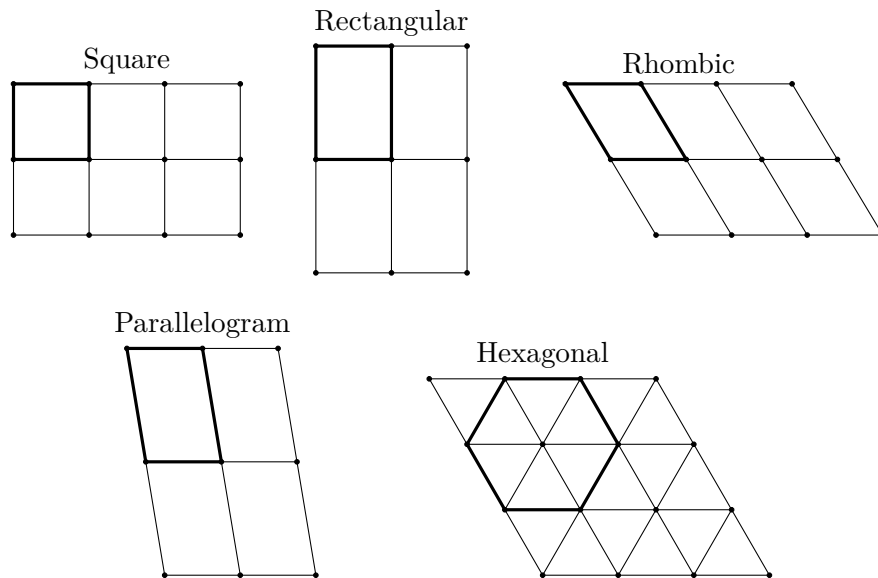
Figure 12.21: Types of lattices in $\mathbb{R}^2$

Notation and Space Groups	Point Group	Lattice Type	Reflections or Glide Reflections?
p1	$\mathbb{Z}_1$	parallelogram	none
p2	$\mathbb{Z}_2$	parallelogram	none
p3	$\mathbb{Z}_3$	hexagonal	none
p4	$\mathbb{Z}_4$	square	none
p6	$\mathbb{Z}_6$	hexagonal	none
pm	D_1	rectangular	reflections
pg	D_1	rectangular	glide reflections
cm	D_1	rhombic	both
pmm	D_2	rectangular	reflections
pmg	D_2	rectangular	glide reflections
pgg	D_2	rectangular	both
c2mm	D_2	rhombic	both
p3m1, p31m	D_3	hexagonal	both
p4m, p4g	D_4	square	both
p6m	D_6	hexagonal	both

Table 12.22: The 17 wallpaper groups
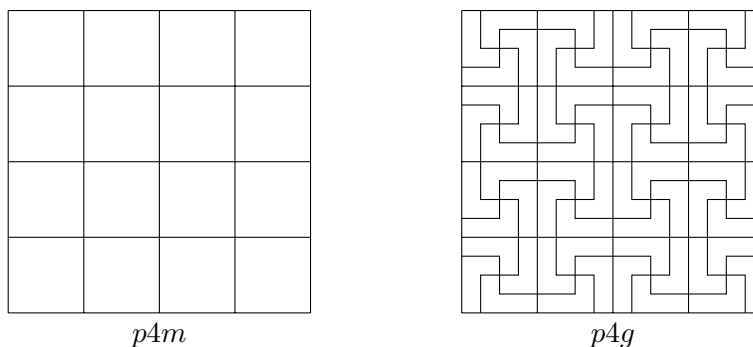
Theorem 12.23. *There are exactly 17 wallpaper groups.*

Figure 12.24: The wallpaper groups p4m and p4g

The 17 wallpaper groups are listed in Table 12.22. The groups p3m1 and p31m can be distinguished by whether or not all of their threefold centers lie on the reflection axes: those of p3m1 must, whereas those of p31m may not. Similarly, the fourfold centers of p4m must lie on the reflection axes whereas those of p4g need not (Figure 12.24). The complete proof of this theorem can be found in several of the references at the end of this chapter, including [5], [6], [10], and [11].

Historical Note

Symmetry groups have intrigued mathematicians for a long time. Leonardo da Vinci was probably the first person to know all of the point groups. At the International Congress of Mathematicians in 1900, David Hilbert gave a now-famous address outlining 23 problems to guide mathematics in the twentieth century. Hilbert's eighteenth problem asked whether or not crystallographic groups in n dimensions were always finite. In 1910, L. Bieberbach proved that crystallographic groups are finite in every dimension. Finding out how many of these groups there are in each dimension is another matter. In $\mathbb{R}^3$ there are 230 different space groups; in $\mathbb{R}^4$ there are 4783. No one has been able to compute the number of space groups for $\mathbb{R}^5$ and beyond. It is interesting to note that the crystallographic groups were found mathematically for $\mathbb{R}^3$ before the 230 different types of crystals were actually discovered in nature.

12.3   Exercises

1.   Prove the identity

$$\langle \mathbf{x}, \mathbf{y} \rangle = \frac{1}{2} \left[\|\mathbf{x} + \mathbf{y}\|^2 - \|\mathbf{x}\|^2 - \|\mathbf{y}\|^2 \right].$$

2.   Show that $O(n)$ is a group.

3.   Prove that the following matrices are orthogonal. Are any of these matrices in $SO(n)$?

(a)

$$\begin{pmatrix} 1/\sqrt{2} & -1/\sqrt{2} \\ 1/\sqrt{2} & 1/\sqrt{2} \end{pmatrix}$$

(b)

$$\begin{pmatrix} 1/\sqrt{5} & 2/\sqrt{5} \\ -2/\sqrt{5} & 1/\sqrt{5} \end{pmatrix}$$

(c)

$$\begin{pmatrix} 4/\sqrt{5} & 0 & 3/\sqrt{5} \\ -3/\sqrt{5} & 0 & 4/\sqrt{5} \\ 0 & -1 & 0 \end{pmatrix}$$

(d)
$$\begin{pmatrix} 1/3 & 2/3 & -2/3 \\ -2/3 & 2/3 & 1/3 \\ -2/3 & 1/3 & 2/3 \end{pmatrix}$$

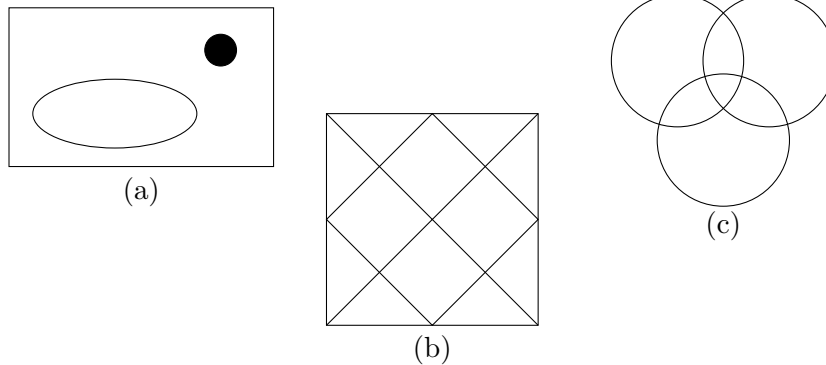**4.** Determine the symmetry group of each of the figures in Figure 12.25.

(a)

(b)

(c)

Figure 12.25

5. Let $\mathbf{x}$, $\mathbf{y}$, and $\mathbf{w}$ be vectors in $\mathbb{R}^n$ and $\alpha \in \mathbb{R}$. Prove each of the following properties of inner products.

(a) $\langle \mathbf{x}, \mathbf{y} \rangle = \langle \mathbf{y}, \mathbf{x} \rangle$.

(b) $\langle \mathbf{x}, \mathbf{y} + \mathbf{w} \rangle = \langle \mathbf{x}, \mathbf{y} \rangle + \langle \mathbf{x}, \mathbf{w} \rangle$.

(c) $\langle \alpha \mathbf{x}, \mathbf{y} \rangle = \langle \mathbf{x}, \alpha \mathbf{y} \rangle = \alpha \langle \mathbf{x}, \mathbf{y} \rangle$.

(d) $\langle \mathbf{x}, \mathbf{x} \rangle \geq 0$ with equality exactly when $\mathbf{x} = 0$.

(e) If $\langle \mathbf{x}, \mathbf{y} \rangle = 0$ for all $\mathbf{x}$ in $\mathbb{R}^n$, then $\mathbf{y} = 0$.

6. Verify that
$$E(n) = \{(A, \mathbf{x}) : A \in O(n) \text{ and } \mathbf{x} \in \mathbb{R}^n\}$$
is a group.

7. Prove that $\{(2, 1), (1, 1)\}$ and $\{(12, 5), (7, 3)\}$ are bases for the same lattice.

8. Let G be a subgroup of $E(2)$ and suppose that T is the translation subgroup of G. Prove that the point group of G is isomorphic to G/T.

9. Let $A \in SL_2(\mathbb{R})$ and suppose that the vectors $\mathbf{x}$ and $\mathbf{y}$ form two sides of a parallelogram in $\mathbb{R}^2$. Prove that the area of this parallelogram is the same as the area of the parallelogram with sides $A\mathbf{x}$ and $A\mathbf{y}$.

10. Prove that $SO(n)$ is a normal subgroup of $O(n)$.

11. Show that any isometry f in $\mathbb{R}^n$ is a one-to-one map.

12. Prove or disprove: an element in $E(2)$ of the form $(A, \mathbf{x})$, where $\mathbf{x} \neq 0$, has infinite order.

13. Prove or disprove: There exists an infinite abelian subgroup of $O(n)$.

14. Let $\mathbf{x} = (x_1, x_2)$ be a point on the unit circle in $\mathbb{R}^2$; that is, $x_1^2 + x_2^2 = 1$. If $A \in O(2)$, show that $A\mathbf{x}$ is also a point on the unit circle.

15. Let G be a group with a subgroup H (not necessarily normal) and a normal subgroup N. Then G is a **semidirect product** of N by H if

- $H \cap N = \{\text{id}\}$;

- $HN = G$.

Show that each of the following is true.

(a) S_3 is the semidirect product of A_3 by $H = \{(1), (12)\}$.

(b) The quaternion group, Q_8, cannot be written as a semidirect product.

(c) $E(2)$ is the semidirect product of $O(2)$ by H, where H consists of all translations in $\mathbb{R}^2$.

16.   Determine which of the 17 wallpaper groups preserves the symmetry of the pattern in Figure 12.16.

17.   Determine which of the 17 wallpaper groups preserves the symmetry of the pattern in Figure 12.26.
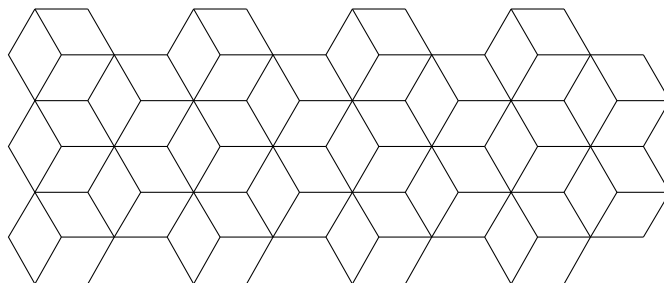
Figure 12.26

18.   Find the rotation group of a dodecahedron.

19.   For each of the 17 wallpaper groups, draw a wallpaper pattern having that group as a symmetry group.

References and Suggested Readings

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[**3**]   Hiller, H. "Crystallography and Cohomology of Groups," *American Mathematical Monthly* **93** (1986), 765–79.

[**4**]   Lockwood, E. H. and Macmillan, R. H. *Geometric Symmetry.* Cambridge University Press, Cambridge, 1978.

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[10] Schattschneider, D. "The Plane Symmetry Groups: Their Recognition and Their Notation," *American Mathematical Monthly* **85**(1978), 439–50.

[11] Schwarzenberger, R. L. "The 17 Plane Symmetry Groups," *Mathematical Gazette* **58**(1974), 123–31.

[12] Weyl, H. *Symmetry.* Princeton University Press, Princeton, NJ, 1952.

12.4 Sage

There is no Sage material for this chapter.

12.5 Sage Exercises

There are no Sage exercises for this chapter.

The Structure of Groups

The ultimate goal of group theory is to classify all groups up to isomorphism; that is, given a particular group, we should be able to match it up with a known group via an isomorphism. For example, we have already proved that any finite cyclic group of order n is isomorphic to $\mathbb{Z}_n$; hence, we "know" all finite cyclic groups. It is probably not reasonable to expect that we will ever know all groups; however, we can often classify certain types of groups or distinguish between groups in special cases.

In this chapter we will characterize all finite abelian groups. We shall also investigate groups with sequences of subgroups. If a group has a sequence of subgroups, say

$$G = H_n \supset H_{n-1} \supset \cdots \supset H_1 \supset H_0 = \{e\},$$

where each subgroup H_i is normal in H_{i+1} and each of the factor groups H_{i+1}/H_i is abelian, then G is a solvable group. In addition to allowing us to distinguish between certain classes of groups, solvable groups turn out to be central to the study of solutions to polynomial equations.

13.1 Finite Abelian Groups

In our investigation of cyclic groups we found that every group of prime order was isomorphic to $\mathbb{Z}_p$, where p was a prime number. We also determined that $\mathbb{Z}_{mn} \cong \mathbb{Z}_m \times \mathbb{Z}_n$ when $\gcd(m, n) = 1$. In fact, much more is true. Every finite abelian group is isomorphic to a direct product of cyclic groups of prime power order; that is, every finite abelian group is isomorphic to a group of the type

$$\mathbb{Z}_{p_1^{\alpha_1}} \times \cdots \times \mathbb{Z}_{p_n^{\alpha_n}},$$

where each p_k is prime (not necessarily distinct).

First, let us examine a slight generalization of finite abelian groups. Suppose that G is a group and let $\{g_i\}$ be a set of elements in G, where i is in some index set I (not necessarily finite). The smallest subgroup of G containing all of the g_i's is the subgroup of G **generated** by the g_i's. If this subgroup of G is in fact all of G, then G is generated by the set $\{g_i : i \in I\}$. In this case the g_i's are said to be the **generators** of G. If there is a finite set $\{g_i : i \in I\}$ that generates G, then G is **finitely generated**.

Example 13.1. Obviously, all finite groups are finitely generated. For example, the group S_3 is generated by the permutations (12) and (123). The group $\mathbb{Z} \times \mathbb{Z}_n$ is an infinite group but is finitely generated by $\{(1, 0), (0, 1)\}$.

Example 13.2. Not all groups are finitely generated. Consider the rational numbers $\mathbb{Q}$ under the operation of addition. Suppose that $\mathbb{Q}$ is finitely generated with generators

$p_1/q_1, \ldots, p_n/q_n$, where each p_i/q_i is a fraction expressed in its lowest terms. Let p be some prime that does not divide any of the denominators $q_1, \ldots, q_n$. We claim that $1/p$ cannot be in the subgroup of $\mathbb{Q}$ that is generated by $p_1/q_1, \ldots, p_n/q_n$, since p does not divide the denominator of any element in this subgroup. This fact is easy to see since the sum of any two generators is

$$p_i/q_i + p_j/q_j = (p_iq_j + p_jq_i)/(q_iq_j).$$

Proposition 13.3. *Let H be the subgroup of a group G that is generated by $\{g_i \in G : i \in I\}$. Then $h \in H$ exactly when it is a product of the form*

$$h = g_{i_1}^{\alpha_1} \cdots g_{i_n}^{\alpha_n},$$

where the g_{i_k}s are not necessarily distinct.

PROOF.   Let K be the set of all products of the form $g_{i_1}^{\alpha_1} \cdots g_{i_n}^{\alpha_n}$, where the g_{i_k}s are not necessarily distinct. Certainly K is a subset of H. We need only show that K is a subgroup of G. If this is the case, then $K = H$, since H is the smallest subgroup containing all the g_is.

Clearly, the set K is closed under the group operation. Since $g_i^0 = 1$, the identity is in K. It remains to show that the inverse of an element $g = g_{i_1}^{k_1} \cdots g_{i_n}^{k_n}$ in K must also be in K. However,

$$g^{-1} = (g_{i_1}^{k_1} \cdots g_{i_n}^{k_n})^{-1} = (g_{i_n}^{-k_n} \cdots g_{i_1}^{-k_1}). \qquad \square$$

The reason that powers of a fixed g_i may occur several times in the product is that we may have a nonabelian group. However, if the group is abelian, then the g_is need occur only once. For example, a product such as $a^{-3}b^5a^7$ in an abelian group could always be simplified (in this case, to a^4b^5).

Now let us restrict our attention to finite abelian groups. We can express any finite abelian group as a finite direct product of cyclic groups. More specifically, letting p be prime, we define a group G to be a *p-**group*** if every element in G has as its order a power of p. For example, both $\mathbb{Z}_2 \times \mathbb{Z}_2$ and $\mathbb{Z}_4$ are 2-groups, whereas $\mathbb{Z}_{27}$ is a 3-group. We shall prove the Fundamental Theorem of Finite Abelian Groups which tells us that every finite abelian group is isomorphic to a direct product of cyclic p-groups.

Theorem 13.4 Fundamental Theorem of Finite Abelian Groups. *Every finite abelian group G is isomorphic to a direct product of cyclic groups of the form*

$$\mathbb{Z}_{p_1^{\alpha_1}} \times \mathbb{Z}_{p_2^{\alpha_2}} \times \cdots \times \mathbb{Z}_{p_n^{\alpha_n}}$$

here the p_i's are primes (not necessarily distinct).

Example 13.5. Suppose that we wish to classify all abelian groups of order $540 = 2^2 \cdot 3^3 \cdot 5$. The Fundamental Theorem of Finite Abelian Groups tells us that we have the following six possibilities.

- $\mathbb{Z}_2 \times \mathbb{Z}_2 \times \mathbb{Z}_3 \times \mathbb{Z}_3 \times \mathbb{Z}_3 \times \mathbb{Z}_5$;

- $\mathbb{Z}_2 \times \mathbb{Z}_2 \times \mathbb{Z}_3 \times \mathbb{Z}_9 \times \mathbb{Z}_5$;

- $\mathbb{Z}_2 \times \mathbb{Z}_2 \times \mathbb{Z}_{27} \times \mathbb{Z}_5$;

- $\mathbb{Z}_4 \times \mathbb{Z}_3 \times \mathbb{Z}_3 \times \mathbb{Z}_3 \times \mathbb{Z}_5$;

- $\mathbb{Z}_4 \times \mathbb{Z}_3 \times \mathbb{Z}_9 \times \mathbb{Z}_5$;

- $\mathbb{Z}_4 \times \mathbb{Z}_{27} \times \mathbb{Z}_5$.

The proof of the Fundamental Theorem of Finite Abelian Groups depends on several lemmas.

Lemma 13.6. *Let G be a finite abelian group of order n. If p is a prime that divides n, then G contains an element of order p.*

PROOF.   We will prove this lemma by induction. If $n = 1$, then there is nothing to show. Now suppose that the order of G is n the lemma is true for all groups of order k, where $k < n$. Furthermore, let p be a prime that divides n.

If G has no proper nontrivial subgroups, then $G = \langle a \rangle$, where a is any element other than the identity. By Exercise 4.4.39, the order of G must be prime. Since p divides n, we know that $p = n$, and G contains $p - 1$ elements of order p.

Now suppose that G contains a nontrivial proper subgroup H. Then $1 < |H| < n$. If $p \mid |H|$, then H contains an element of order p by induction and the lemma is true. Suppose that p does not divide the order of H. Since G is abelian, it must be the case that H is a normal subgroup of G, and $|G| = |H| \cdot |G/H|$. Consequently, p must divide $|G/H|$. Since $|G/H| < |G| = n$, we know that G/H contains an element aH of order p by the induction hypothesis. Thus,

$$H = (aH)^p = a^p H,$$

and $a^p \in H$ but $a \notin H$. If $|H| = r$, then p and r are relatively prime, and there exist integers s and t such that $sp + tr = 1$. Furthermore, the order of a^p must divide r, and $(a^p)^r = (a^r)^p = 1$.

We claim that a^r has order p. We must show that $a^r \neq 1$. Suppose $a^r = 1$. Then

$$
\begin{aligned}
a &= a^{sp+tr} \\
  &= a^{sp} a^{tr} \\
  &= (a^p)^s (a^r)^t \\
  &= (a^p)^s 1 \\
  &= (a^p)^s.
\end{aligned}
$$

Since $a^p \in H$, it must be the case that $a = (a^p)^s \in H$, which is a contradiction. Therefore, $a^r \neq 1$ is an element of order p in G. $\square$

Lemma 13.6 is a special case of Cauchy's Theorem (Theorem 15.1), which states that if G be a finite group and p a prime such that p divides the order of G, then G contains a subgroup of order p. We will prove Cauchy's Theorem in Chapter 15.

Lemma 13.7. *A finite abelian group is a p-group if and only if its order is a power of p.*

PROOF.   If $|G| = p^n$ then by Lagrange's theorem, then the order of any $g \in G$ must divide p^n, and therefore must be a power of p. Conversely, if $|G|$ is not a power of p, then it has some other prime divisor q, so by Lemma 13.6, G has an element of order q and thus is not a p-group. $\square$

Lemma 13.8. *Let G be a finite abelian group of order $n = p_1^{\alpha_1} \cdots p_k^{\alpha_k}$, where where $p_1, \ldots, p_k$ are distinct primes and $\alpha_1, \alpha_2, \ldots, \alpha_k$ are positive integers. Then G is the internal direct product of subgroups $G_1, G_2, \ldots, G_k$, where G_i is the subgroup of G consisting of all elements of order p_i^k for some integer k.*

PROOF. Since G is an abelian group, we are guaranteed that G_i is a subgroup of G for $i = 1, \ldots, n$. Since the identity has order $p_i^0 = 1$, we know that $1 \in G_i$. If $g \in G_i$ has order p_i^r, then g^{-1} must also have order p_i^r. Finally, if $h \in G_i$ has order p_i^s, then

$$(gh)^{p_i^t} = g^{p_i^t} h^{p_i^t} = 1 \cdot 1 = 1,$$

where t is the maximum of r and s.

We must show that

$$G = G_1 G_2 \cdots G_k$$

and $G_i \cap G_j = \{1\}$ for $i \neq j$. Suppose that $g_1 \in G_1$ is in the subgroup generated by $G_2, G_3, \ldots, G_k$. Then $g_1 = g_2 g_3 \cdots g_k$ for $g_i \in G_i$. Since g_i has order p^{α_i}, we know that $g_i^{p^{\alpha_i}} = 1$ for $i = 2, 3, \ldots, k$, and $g_1^{p_2^{\alpha_2} \cdots p_k^{\alpha_k}} = 1$. Since the order of g_1 is a power of p_1 and $\gcd(p_1, p_2^{\alpha_2} \cdots p_k^{\alpha_k}) = 1$, it must be the case that $g_1 = 1$ and the intersection of G_1 with any of the subgroups $G_2, G_3, \ldots, G_k$ is the identity. A similar argument shows that $G_i \cap G_j = \{1\}$ for $i \neq j$.

Next, we must show that it possible to write every $g \in G$ as a product $g_1 \cdots g_k$, where $g_i \in G_i$. Since the order of g divides the order of G, we know that

$$|g| = p_1^{\beta_1} p_2^{\beta_2} \cdots p_k^{\beta_k}$$

for some integers $\beta_1, \ldots, \beta_k$. Letting $a_i = |g|/p_i^{\beta_i}$, the a_i's are relatively prime; hence, there exist integers $b_1, \ldots, b_k$ such that $a_1 b_1 + \cdots + a_k b_k = 1$. Consequently,

$$g = g^{a_1 b_1 + \cdots + a_k b_k} = g^{a_1 b_1} \cdots g^{a_k b_k}.$$

Since

$$g^{(a_i b_i) p_i^{\beta_i}} = g^{b_i |g|} = e,$$

it follows that $g^{a_i b_i}$ must be in G_i. Let $g_i = g^{a_i b_i}$. Then $g = g_1 \cdots g_k \in G_1 G_2 \cdots G_k$. Therefore, $G = G_1 G_2 \cdots G_k$ is an internal direct product of subgroups. $\square$

If remains for us to determine the possible structure of each p_i-group G_i in Lemma 13.8.

Lemma 13.9. *Let G be a finite abelian p-group and suppose that $g \in G$ has maximal order. Then G is isomorphic to $\langle g \rangle \times H$ for some subgroup H of G.*

PROOF. By Lemma 13.7, we may assume that the order of G is p^n. We shall induct on n. If $n = 1$, then G is cyclic of order p and must be generated by g. Suppose now that the statement of the lemma holds for all integers k with $1 \leq k < n$ and let g be of maximal order in G, say $|g| = p^m$. Then $a^{p^m} = e$ for all $a \in G$. Now choose h in G such that $h \notin \langle g \rangle$, where h has the smallest possible order. Certainly such an h exists; otherwise, $G = \langle g \rangle$ and we are done. Let $H = \langle h \rangle$.

We claim that $\langle g \rangle \cap H = \{e\}$. It suffices to show that $|H| = p$. Since $|h^p| = |h|/p$, the order of h^p is smaller than the order of h and must be in $\langle g \rangle$ by the minimality of h; that is, $h^p = g^r$ for some number r. Hence,

$$(g^r)^{p^{m-1}} = (h^p)^{p^{m-1}} = h^{p^m} = e,$$

and the order of g^r must be less than or equal to p^{m-1}. Therefore, g^r cannot generate $\langle g \rangle$. Notice that p must occur as a factor of r, say $r = ps$, and $h^p = g^r = g^{ps}$. Define a to be $g^{-s} h$. Then a cannot be in $\langle g \rangle$; otherwise, h would also have to be in $\langle g \rangle$. Also,

$$a^p = g^{-sp} h^p = g^{-r} h^p = h^{-p} h^p = e.$$

We have now formed an element a with order p such that $a \notin \langle g \rangle$. Since h was chosen to have the smallest order of all of the elements that are not in $\langle g \rangle$, $|H| = p$.

Now we will show that the order of gH in the factor group G/H must be the same as the order of g in G. If $|gH| < |g| = p^m$, then

$$H = (gH)^{p^{m-1}} = g^{p^{m-1}} H;$$

hence, $g^{p^{m-1}}$ must be in $\langle g \rangle \cap H = \{e\}$, which contradicts the fact that the order of g is p^m. Therefore, gH must have maximal order in G/H. By the Correspondence Theorem and our induction hypothesis,

$$G/H \cong \langle gH \rangle \times K/H$$

for some subgroup K of G containing H. We claim that $\langle g \rangle \cap K = \{e\}$. If $b \in \langle g \rangle \cap K$, then $bH \in \langle gH \rangle \cap K/H = \{H\}$ and $b \in \langle g \rangle \cap H = \{e\}$. It follows that $G = \langle g \rangle K$ implies that $G \cong \langle g \rangle \times K$. $\qquad\square$

The proof of the Fundamental Theorem of Finite Abelian Groups follows very quickly from Lemma 13.9. Suppose that G is a finite abelian group and let g be an element of maximal order in G. If $\langle g \rangle = G$, then we are done; otherwise, $G \cong \mathbb{Z}_{|g|} \times H$ for some subgroup H contained in G by the lemma. Since $|H| < |G|$, we can apply mathematical induction.

We now state the more general theorem for all finitely generated abelian groups. The proof of this theorem can be found in any of the references at the end of this chapter.

Theorem 13.10 The Fundamental Theorem of Finitely Generated Abelian Groups.
Every finitely generated abelian group G is isomorphic to a direct product of cyclic groups of the form

$$\mathbb{Z}_{p_1^{\alpha_1}} \times \mathbb{Z}_{p_2^{\alpha_2}} \times \cdots \times \mathbb{Z}_{p_n^{\alpha_n}} \times \mathbb{Z} \times \cdots \times \mathbb{Z},$$

where the p_i's are primes (not necessarily distinct).

13.2   Solvable Groups

A **subnormal series** of a group G is a finite sequence of subgroups

$$G = H_n \supset H_{n-1} \supset \cdots \supset H_1 \supset H_0 = \{e\},$$

where H_i is a normal subgroup of H_{i+1}. If each subgroup H_i is normal in G, then the series is called a **normal series**. The **length** of a subnormal or normal series is the number of proper inclusions.

Example 13.11. Any series of subgroups of an abelian group is a normal series. Consider the following series of groups:

$$\mathbb{Z} \supset 9\mathbb{Z} \supset 45\mathbb{Z} \supset 180\mathbb{Z} \supset \{0\},$$
$$\mathbb{Z}_{24} \supset \langle 2 \rangle \supset \langle 6 \rangle \supset \langle 12 \rangle \supset \{0\}.$$

Example 13.12. A subnormal series need not be a normal series. Consider the following subnormal series of the group D_4:

$$D_4 \supset \{(1), (12)(34), (13)(24), (14)(23)\} \supset \{(1), (12)(34)\} \supset \{(1)\}.$$

The subgroup $\{(1), (12)(34)\}$ is not normal in D_4; consequently, this series is not a normal series.

A subnormal (normal) series $\{K_j\}$ is a **refinement of a subnormal (normal) series** $\{H_i\}$ if $\{H_i\} \subset \{K_j\}$. That is, each H_i is one of the K_j.

Example 13.13. The series

$$\mathbb{Z} \supset 3\mathbb{Z} \supset 9\mathbb{Z} \supset 45\mathbb{Z} \supset 90\mathbb{Z} \supset 180\mathbb{Z} \supset \{0\}$$

is a refinement of the series

$$\mathbb{Z} \supset 9\mathbb{Z} \supset 45\mathbb{Z} \supset 180\mathbb{Z} \supset \{0\}.$$

The best way to study a subnormal or normal series of subgroups, $\{H_i\}$ of G, is actually to study the factor groups H_{i+1}/H_i. We say that two subnormal (normal) series $\{H_i\}$ and $\{K_j\}$ of a group G are **isomorphic** if there is a one-to-one correspondence between the collections of factor groups $\{H_{i+1}/H_i\}$ and $\{K_{j+1}/K_j\}$.

Example 13.14. The two normal series

$$\mathbb{Z}_{60} \supset \langle 3 \rangle \supset \langle 15 \rangle \supset \{0\}$$
$$\mathbb{Z}_{60} \supset \langle 4 \rangle \supset \langle 20 \rangle \supset \{0\}$$

of the group $\mathbb{Z}_{60}$ are isomorphic since

$$\mathbb{Z}_{60}/\langle 3 \rangle \cong \langle 20 \rangle/\{0\} \cong \mathbb{Z}_3$$
$$\langle 3 \rangle/\langle 15 \rangle \cong \langle 4 \rangle/\langle 20 \rangle \cong \mathbb{Z}_5$$
$$\langle 15 \rangle/\{0\} \cong \mathbb{Z}_{60}/\langle 4 \rangle \cong \mathbb{Z}_4.$$

A subnormal series $\{H_i\}$ of a group G is a **composition series** if all the factor groups are simple; that is, if none of the factor groups of the series contains a normal subgroup. A normal series $\{H_i\}$ of G is a **principal series** if all the factor groups are simple.

Example 13.15. The group $\mathbb{Z}_{60}$ has a composition series

$$\mathbb{Z}_{60} \supset \langle 3 \rangle \supset \langle 15 \rangle \supset \langle 30 \rangle \supset \{0\}$$

with factor groups

$$\mathbb{Z}_{60}/\langle 3 \rangle \cong \mathbb{Z}_3$$
$$\langle 3 \rangle/\langle 15 \rangle \cong \mathbb{Z}_5$$
$$\langle 15 \rangle/\langle 30 \rangle \cong \mathbb{Z}_2$$
$$\langle 30 \rangle/\{0\} \cong \mathbb{Z}_2.$$

Since $\mathbb{Z}_{60}$ is an abelian group, this series is automatically a principal series. Notice that a composition series need not be unique. The series

$$\mathbb{Z}_{60} \supset \langle 2 \rangle \supset \langle 4 \rangle \supset \langle 20 \rangle \supset \{0\}$$

is also a composition series.

Example 13.16. For $n \geq 5$, the series

$$S_n \supset A_n \supset \{(1)\}$$

is a composition series for S_n since $S_n/A_n \cong \mathbb{Z}_2$ and A_n is simple.

Example 13.17. Not every group has a composition series or a principal series. Suppose that

$$\{0\} = H_0 \subset H_1 \subset \cdots \subset H_{n-1} \subset H_n = \mathbb{Z}$$

is a subnormal series for the integers under addition. Then H_1 must be of the form $k\mathbb{Z}$ for some $k \in \mathbb{N}$. In this case $H_1/H_0 \cong k\mathbb{Z}$ is an infinite cyclic group with many nontrivial proper normal subgroups.

Although composition series need not be unique as in the case of $\mathbb{Z}_{60}$, it turns out that any two composition series are related. The factor groups of the two composition series for $\mathbb{Z}_{60}$ are $\mathbb{Z}_2$, $\mathbb{Z}_2$, $\mathbb{Z}_3$, and $\mathbb{Z}_5$; that is, the two composition series are isomorphic. The Jordan-Hölder Theorem says that this is always the case.

Theorem 13.18 Jordan-Hölder. *Any two composition series of G are isomorphic.*

PROOF.   We shall employ mathematical induction on the length of the composition series. If the length of a composition series is 1, then G must be a simple group. In this case any two composition series are isomorphic.

Suppose now that the theorem is true for all groups having a composition series of length k, where $1 \leq k < n$. Let

$$G = H_n \supset H_{n-1} \supset \cdots \supset H_1 \supset H_0 = \{e\}$$
$$G = K_m \supset K_{m-1} \supset \cdots \supset K_1 \supset K_0 = \{e\}$$

be two composition series for G. We can form two new subnormal series for G since $H_i \cap K_{m-1}$ is normal in $H_{i+1} \cap K_{m-1}$ and $K_j \cap H_{n-1}$ is normal in $K_{j+1} \cap H_{n-1}$:

$$G = H_n \supset H_{n-1} \supset H_{n-1} \cap K_{m-1} \supset \cdots \supset H_0 \cap K_{m-1} = \{e\}$$
$$G = K_m \supset K_{m-1} \supset K_{m-1} \cap H_{n-1} \supset \cdots \supset K_0 \cap H_{n-1} = \{e\}.$$

Since $H_i \cap K_{m-1}$ is normal in $H_{i+1} \cap K_{m-1}$, the Second Isomorphism Theorem (Theorem 11.12) implies that

$$(H_{i+1} \cap K_{m-1})/(H_i \cap K_{m-1}) = (H_{i+1} \cap K_{m-1})/(H_i \cap (H_{i+1} \cap K_{m-1}))$$
$$\cong H_i(H_{i+1} \cap K_{m-1})/H_i,$$

where H_i is normal in $H_i(H_{i+1} \cap K_{m-1})$. Since $\{H_i\}$ is a composition series, H_{i+1}/H_i must be simple; consequently, $H_i(H_{i+1} \cap K_{m-1})/H_i$ is either H_{i+1}/H_i or H_i/H_i. That is, $H_i(H_{i+1} \cap K_{m-1})$ must be either H_i or H_{i+1}. Removing any nonproper inclusions from the series

$$H_{n-1} \supset H_{n-1} \cap K_{m-1} \supset \cdots \supset H_0 \cap K_{m-1} = \{e\},$$

we have a composition series for H_{n-1}. Our induction hypothesis says that this series must be equivalent to the composition series

$$H_{n-1} \supset \cdots \supset H_1 \supset H_0 = \{e\}.$$

Hence, the composition series

$$G = H_n \supset H_{n-1} \supset \cdots \supset H_1 \supset H_0 = \{e\}$$

and

$$G = H_n \supset H_{n-1} \supset H_{n-1} \cap K_{m-1} \supset \cdots \supset H_0 \cap K_{m-1} = \{e\}$$

are equivalent. If $H_{n-1} = K_{m-1}$, then the composition series $\{H_i\}$ and $\{K_j\}$ are equivalent and we are done; otherwise, $H_{n-1}K_{m-1}$ is a normal subgroup of G properly containing

H_{n-1}. In this case $H_{n-1}K_{m-1} = G$ and we can apply the Second Isomorphism Theorem once again; that is,

$$K_{m-1}/(K_{m-1} \cap H_{n-1}) \cong (H_{n-1}K_{m-1})/H_{n-1} = G/H_{n-1}.$$

Therefore,

$$G = H_n \supset H_{n-1} \supset H_{n-1} \cap K_{m-1} \supset \cdots \supset H_0 \cap K_{m-1} = \{e\}$$

and

$$G = K_m \supset K_{m-1} \supset K_{m-1} \cap H_{n-1} \supset \cdots \supset K_0 \cap H_{n-1} = \{e\}$$

are equivalent and the proof of the theorem is complete. $\square$

A group G is **solvable** if it has a subnormal series $\{H_i\}$ such that all of the factor groups H_{i+1}/H_i are abelian. Solvable groups will play a fundamental role when we study Galois theory and the solution of polynomial equations.

Example 13.19. The group S_4 is solvable since

$$S_4 \supset A_4 \supset \{(1), (12)(34), (13)(24), (14)(23)\} \supset \{(1)\}$$

has abelian factor groups; however, for $n \geq 5$ the series

$$S_n \supset A_n \supset \{(1)\}$$

is a composition series for S_n with a nonabelian factor group. Therefore, S_n is not a solvable group for $n \geq 5$.

13.3   Exercises

1.   Find all of the abelian groups of order less than or equal to 40 up to isomorphism.

2.   Find all of the abelian groups of order 200 up to isomorphism.

3.   Find all of the abelian groups of order 720 up to isomorphism.

4.   Find all of the composition series for each of the following groups.

(a) $\mathbb{Z}_{12}$

(b) $\mathbb{Z}_{48}$

(c) The quaternions, Q_8

(d) D_4

(e) $S_3 \times \mathbb{Z}_4$

(f) S_4

(g) $S_n, n \geq 5$

(h) $\mathbb{Q}$

5.   Show that the infinite direct product $G = \mathbb{Z}_2 \times \mathbb{Z}_2 \times \cdots$ is not finitely generated.

6.   Let G be an abelian group of order m. If n divides m, prove that G has a subgroup of order n.

7.   A group G is a **torsion group** if every element of G has finite order. Prove that a finitely generated abelian torsion group must be finite.

8.   Let G, H, and K be finitely generated abelian groups. Show that if $G \times H \cong G \times K$, then $H \cong K$. Give a counterexample to show that this cannot be true in general.

9.   Let G and H be solvable groups. Show that $G \times H$ is also solvable.

10.   If G has a composition (principal) series and if N is a proper normal subgroup of G, show there exists a composition (principal) series containing N.

11.   Prove or disprove: Let N be a normal subgroup of G. If N and G/N have composition series, then G must also have a composition series.

12.   Let N be a normal subgroup of G. If N and G/N are solvable groups, show that G is also a solvable group.

13.   Prove that G is a solvable group if and only if G has a series of subgroups

$$G = P_n \supset P_{n-1} \supset \cdots \supset P_1 \supset P_0 = \{e\}$$

where P_i is normal in P_{i+1} and the order of P_{i+1}/P_i is prime.

14.   Let G be a solvable group. Prove that any subgroup of G is also solvable.

15.   Let G be a solvable group and N a normal subgroup of G. Prove that G/N is solvable.

16.   Prove that D_n is solvable for all integers n.

17.   Suppose that G has a composition series. If N is a normal subgroup of G, show that N and G/N also have composition series.

18.   Let G be a cyclic p-group with subgroups H and K. Prove that either H is contained in K or K is contained in H.

19.   Suppose that G is a solvable group with order $n \geq 2$. Show that G contains a normal nontrivial abelian subgroup.

20.   Recall that the ***commutator subgroup*** G' of a group G is defined as the subgroup of G generated by elements of the form $a^{-1}b^{-1}ab$ for $a, b \in G$. We can define a series of subgroups of G by $G^{(0)} = G$, $G^{(1)} = G'$, and $G^{(i+1)} = (G^{(i)})'$.

(a) Prove that $G^{(i+1)}$ is normal in $(G^{(i)})'$. The series of subgroups

$$G^{(0)} = G \supset G^{(1)} \supset G^{(2)} \supset \cdots$$

is called the ***derived series*** of G.

(b) Show that G is solvable if and only if $G^{(n)} = \{e\}$ for some integer n.

21.   Suppose that G is a solvable group with order $n \geq 2$. Show that G contains a normal nontrivial abelian factor group.

22. **Zassenhaus Lemma.**   Let H and K be subgroups of a group G. Suppose also that H^* and K^* are normal subgroups of H and K respectively. Then

(a) $H^*(H \cap K^*)$ is a normal subgroup of $H^*(H \cap K)$.

(b) $K^*(H^* \cap K)$ is a normal subgroup of $K^*(H \cap K)$.

(c) $H^*(H \cap K)/H^*(H \cap K^*) \cong K^*(H \cap K)/K^*(H^* \cap K) \cong (H \cap K)/(H^* \cap K)(H \cap K^*)$.

23. **Schreier's Theorem.**   Use the Zassenhaus Lemma to prove that two subnormal (normal) series of a group G have isomorphic refinements.

24.   Use Schreier's Theorem to prove the Jordan-Hölder Theorem.

13.4   Programming Exercises

1.   Write a program that will compute all possible abelian groups of order n. What is the largest n for which your program will work?

References and Suggested Readings

[**1**]   Hungerford, T. W. *Algebra.* Springer, New York, 1974.

[**2**]   Lang, S. *Algebra.* 3rd ed. Springer, New York, 2002.

[**3**]   Rotman, J. J. *An Introduction to the Theory of Groups.* 4th ed. Springer, New York, 1995.

13.5   Sage

Cyclic groups, and direct products of cyclic groups, are implemented in Sage as permutation groups. However, these groups quickly become very unwieldly representations and it should be easier to work with finite abelian groups in Sage. So we will postpone any specifics for this chapter until that happens. However, now that we understand the notion of isomorphic groups and the structure of finite abelian groups, we can return to our quest to classify all of the groups with order less than 16.

Classification of Finite Groups

It does not take any sophisticated tools to understand groups of order $2p$, where p is an odd prime. There are two possibilities — a cyclic group of order $2p$ and the dihedral group of order $2p$ that is the set of symmetries of a regular p-gon. The proof requires some close, tight reasoning, but the required theorems are generally just concern orders of elements, Lagrange's Theorem and cosets. See Exercise 9.3.55. This takes care of orders $n = 6$, 10, 14.

For $n = 9$, the upcoming Corollary 14.16 will tell us that any group of order p^2 (where p is a prime) is abelian. So we know from this section that the only two possibilities are $\mathbb{Z}_9$ and $\mathbb{Z}_3 \times \mathbb{Z}_3$. Similarly, the upcoming Theorem 15.10 will tell us that every group of order $n = 15$ is abelian. Now this leaves just one possibility for this order: $\mathbb{Z}_3 \times \mathbb{Z}_5 \cong \mathbb{Z}_{15}$.

We have just two orders left to analyze: $n = 8$ and $n = 12$. The possibilities are groups we already know, with one exception. However, the analysis that these are the *only* possibilities is more complicated, and will not be pursued now, nor in the next few chapters. Notice that $n = 16$ is more complicated still, with 14 different possibilities (which explains why we stopped here).

For $n = 8$ there are 3 abelian groups, and the two non-abelian groups are the dihedral group (symmetries of a square) and the quaternions.

For $n = 12$ there are 2 abelian groups, and 3 non-abelian groups. We know two of the non-abelian groups as a dihedral group, and the alternating group on 4 symbols (which is also the symmetries of a tetrahedron). The third non-abelian group is an example of a "dicyclic" group, which is an infinite family of groups, each with order divisible by 4. The order 12 dicyclic group can also be constructed as a "semi-direct product" of two cyclic groups — this is a construction worth knowing as you pursue further study of group theory. The order 8 dicyclic group is also the quaternions and more generally, the dicyclic groups of order 2^k, $k > 2$ are known as "generalized quaternion groups."

The following examples will show you how to construct some of these groups, while also exercising a few of the commands and allowing us to be more certain the following table is accurate.

```
S = SymmetricGroup(3)
D = DihedralGroup(3)
S.is_isomorphic(D)
```

```
True
```

```
C3 = CyclicPermutationGroup(3)
C5 = CyclicPermutationGroup(5)
DP = direct_product_permgroups([C3, C5])
C  = CyclicPermutationGroup(15)
DP.is_isomorphic(C)
```

```
True
```

```
Q  = QuaternionGroup()
DI = DiCyclicGroup(2)
Q.is_isomorphic(DI)
```

```
True
```

Groups of Small Order as Permutation Groups

We list here constructions, as permutation groups in Sage, for all of the groups of order less than 16.

Order	Construction	Notes, Alternatives
1	CyclicPermutationGroup(1)	Trivial
2	CyclicPermutationGroup(2)	SymmetricGroup(2)
3	CyclicPermutationGroup(3)	Prime order
4	CyclicPermutationGroup(4)	Cyclic
4	KleinFourGroup()	Abelian, non-cyclic
5	CyclicPermutationGroup(5)	Prime order
6	CyclicPermutationGroup(6)	Cyclic
6	SymmetricGroup(3)	Non-abelian DihedralGroup(3)
7	CyclicPermutationGroup(7)	Prime order
8	CyclicPermutationGroup(8)	Cyclic
8	C2=CyclicPermutationGroup(2) C4=CyclicPermutationGroup(4) G=direct_product_permgroups([C2,C4])	Abelian, non-cyclic
8	C2=CyclicPermutationGroup(2) G=direct_product_permgroups([C2,C2,C2])	Abelian, non-cyclic
8	DihedralGroup(4)	Non-abelian
8	QuaternionGroup()	Quaternions DiCyclicGroup(2)
9	CyclicPermutationGroup(9)	Cyclic
9	C3=CyclicPermutationGroup(3) G=direct_product_permgroups([C3,C3])	Abelian, non-cyclic
10	CyclicPermutationGroup(10)	Cyclic
10	DihedralGroup(5)	Non-abelian
11	CyclicPermutationGroup(11)	Prime order
12	CyclicPermutationGroup(12)	Cyclic
12	C2=CyclicPermutationGroup(2) C6=CyclicPermutationGroup(6) G=direct_product_permgroups([C2,C6])	Abelian, non-cyclic
12	DihedralGroup(6)	Non-abelian
12	AlternatingGroup(4)	Non-abelian Symmetries of tetrahedron
12	DiCyclicGroup(3)	Non-abelian Semi-direct product $Z_3 \rtimes Z_4$
13	CyclicPermutationGroup(13)	Prime order
14	CyclicPermutationGroup(14)	Cyclic
14	DihedralGroup(7)	Non-abelian
15	CyclicPermutationGroup(15)	Cyclic

Table 13.20: The Groups of Order **16** or Less in Sage

13.6   Sage Exercises

There are no Sage exercises for this chapter.

14

Group Actions

Group actions generalize group multiplication. If G is a group and X is an arbitrary set, a group action of an element $g \in G$ and $x \in X$ is a product, gx, living in X. Many problems in algebra are best be attacked via group actions. For example, the proofs of the Sylow theorems and of Burnside's Counting Theorem are most easily understood when they are formulated in terms of group actions.

14.1 Groups Acting on Sets

Let X be a set and G be a group. A *(left) action* of G on X is a map $G \times X \to X$ given by $(g, x) \mapsto gx$, where

1. $ex = x$ for all $x \in X$;

2. $(g_1 g_2)x = g_1(g_2 x)$ for all $x \in X$ and all $g_1, g_2 \in G$.

Under these considerations X is called a *G-set*. Notice that we are not requiring X to be related to G in any way. It is true that every group G acts on every set X by the trivial action $(g, x) \mapsto x$; however, group actions are more interesting if the set X is somehow related to the group G.

Example 14.1. Let $G = GL_2(\mathbb{R})$ and $X = \mathbb{R}^2$. Then G acts on X by left multiplication. If $v \in \mathbb{R}^2$ and I is the identity matrix, then $Iv = v$. If A and B are 2×2 invertible matrices, then $(AB)v = A(Bv)$ since matrix multiplication is associative.

Example 14.2. Let $G = D_4$ be the symmetry group of a square. If $X = \{1, 2, 3, 4\}$ is the set of vertices of the square, then we can consider D_4 to consist of the following permutations:

$$\{(1), (13), (24), (1432), (1234), (12)(34), (14)(23), (13)(24)\}.$$

The elements of D_4 act on X as functions. The permutation $(13)(24)$ acts on vertex 1 by sending it to vertex 3, on vertex 2 by sending it to vertex 4, and so on. It is easy to see that the axioms of a group action are satisfied.

In general, if X is any set and G is a subgroup of S_X, the group of all permutations acting on X, then X is a G-set under the group action

$$(\sigma, x) \mapsto \sigma(x)$$

for $\sigma \in G$ and $x \in X$.

Example 14.3. If we let $X = G$, then every group G acts on itself by the left regular representation; that is, $(g, x) \mapsto \lambda_g(x) = gx$, where λ_g is left multiplication:

$$e \cdot x = \lambda_e x = ex = x$$
$$(gh) \cdot x = \lambda_{gh} x = \lambda_g \lambda_h x = \lambda_g(hx) = g \cdot (h \cdot x).$$

If H is a subgroup of G, then G is an H-set under left multiplication by elements of H.

Example 14.4. Let G be a group and suppose that $X = G$. If H is a subgroup of G, then G is an H-set under **conjugation**; that is, we can define an action of H on G,

$$H \times G \to G,$$

via

$$(h, g) \mapsto hgh^{-1}$$

for $h \in H$ and $g \in G$. Clearly, the first axiom for a group action holds. Observing that

$$\begin{aligned} (h_1 h_2, g) &= h_1 h_2 g (h_1 h_2)^{-1} \\ &= h_1 (h_2 g h_2^{-1}) h_1^{-1} \\ &= (h_1, (h_2, g)), \end{aligned}$$

we see that the second condition is also satisfied.

Example 14.5. Let H be a subgroup of G and $\mathcal{L}_H$ the set of left cosets of H. The set $\mathcal{L}_H$ is a G-set under the action

$$(g, xH) \mapsto gxH.$$

Again, it is easy to see that the first axiom is true. Since $(gg')xH = g(g'xH)$, the second axiom is also true.

If G acts on a set X and $x, y \in X$, then x is said to be *G-equivalent* to y if there exists a $g \in G$ such that $gx = y$. We write $x \sim_G y$ or $x \sim y$ if two elements are G-equivalent.

Proposition 14.6. *Let X be a G-set. Then G-equivalence is an equivalence relation on X.*

PROOF. The relation $\sim$ is reflexive since $ex = x$. Suppose that $x \sim y$ for $x, y \in X$. Then there exists a g such that $gx = y$. In this case $g^{-1}y = x$; hence, $y \sim x$. To show that the relation is transitive, suppose that $x \sim y$ and $y \sim z$. Then there must exist group elements g and h such that $gx = y$ and $hy = z$. So $z = hy = (hg)x$, and x is equivalent to z. $\square$

If X is a G-set, then each partition of X associated with G-equivalence is called an **orbit** of X under G. We will denote the orbit that contains an element x of X by $\mathcal{O}_x$.

Example 14.7. Let G be the permutation group defined by

$$G = \{(1), (123), (132), (45), (123)(45), (132)(45)\}$$

and $X = \{1, 2, 3, 4, 5\}$. Then X is a G-set. The orbits are $\mathcal{O}_1 = \mathcal{O}_2 = \mathcal{O}_3 = \{1, 2, 3\}$ and $\mathcal{O}_4 = \mathcal{O}_5 = \{4, 5\}$.

Now suppose that G is a group acting on a set X and let g be an element of G. The **fixed point set** of g in X, denoted by X_g, is the set of all $x \in X$ such that $gx = x$. We can also study the group elements g that fix a given $x \in X$. This set is more than a subset of G, it is a subgroup. This subgroup is called the **stabilizer subgroup** or **isotropy subgroup** of x. We will denote the stabilizer subgroup of x by G_x.

Remark 14.8. It is important to remember that $X_g \subset X$ and $G_x \subset G$.

Example 14.9. Let $X = \{1, 2, 3, 4, 5, 6\}$ and suppose that G is the permutation group given by the permutations

$$\{(1), (12)(3456), (35)(46), (12)(3654)\}.$$

Then the fixed point sets of X under the action of G are

$$X_{(1)} = X,$$
$$X_{(35)(46)} = \{1, 2\},$$
$$X_{(12)(3456)} = X_{(12)(3654)} = \emptyset,$$

and the stabilizer subgroups are

$$G_1 = G_2 = \{(1), (35)(46)\},$$
$$G_3 = G_4 = G_5 = G_6 = \{(1)\}.$$

It is easily seen that G_x is a subgroup of G for each $x \in X$.

Proposition 14.10. *Let G be a group acting on a set X and $x \in X$. The stabilizer group of x, G_x, is a subgroup of G.*

PROOF.   Clearly, $e \in G_x$ since the identity fixes every element in the set X. Let $g, h \in G_x$. Then $gx = x$ and $hx = x$. So $(gh)x = g(hx) = gx = x$; hence, the product of two elements in G_x is also in G_x. Finally, if $g \in G_x$, then $x = ex = (g^{-1}g)x = (g^{-1})gx = g^{-1}x$. So g^{-1} is in G_x.   □

We will denote the number of elements in the fixed point set of an element $g \in G$ by $|X_g|$ and denote the number of elements in the orbit of $x \in X$ by $|\mathcal{O}_x|$. The next theorem demonstrates the relationship between orbits of an element $x \in X$ and the left cosets of G_x in G.

Theorem 14.11. *Let G be a finite group and X a finite G-set. If $x \in X$, then $|\mathcal{O}_x| = [G : G_x]$.*

PROOF.   We know that $|G|/|G_x|$ is the number of left cosets of G_x in G by Lagrange's Theorem (Theorem 6.10). We will define a bijective map ϕ between the orbit $\mathcal{O}_x$ of X and the set of left cosets $\mathcal{L}_{G_x}$ of G_x in G. Let $y \in \mathcal{O}_x$. Then there exists a g in G such that $gx = y$. Define ϕ by $\phi(y) = gG_x$. To show that ϕ is one-to-one, assume that $\phi(y_1) = \phi(y_2)$. Then

$$\phi(y_1) = g_1 G_x = g_2 G_x = \phi(y_2),$$

where $g_1 x = y_1$ and $g_2 x = y_2$. Since $g_1 G_x = g_2 G_x$, there exists a $g \in G_x$ such that $g_2 = g_1 g$,

$$y_2 = g_2 x = g_1 g x = g_1 x = y_1;$$

consequently, the map ϕ is one-to-one. Finally, we must show that the map ϕ is onto. Let gG_x be a left coset. If $gx = y$, then $\phi(y) = gG_x$.   □

14.2 The Class Equation

Let X be a finite G-set and X_G be the set of fixed points in X; that is,

$$X_G = \{x \in X : gx = x \text{ for all } g \in G\}.$$

Since the orbits of the action partition X,

$$|X| = |X_G| + \sum_{i=k}^{n} |\mathcal{O}_{x_i}|,$$

where $x_k, \ldots, x_n$ are representatives from the distinct nontrivial orbits of X.

Now consider the special case in which G acts on itself by conjugation, $(g, x) \mapsto gxg^{-1}$. The **center** of G,

$$Z(G) = \{x : xg = gx \text{ for all } g \in G\},$$

is the set of points that are fixed by conjugation. The nontrivial orbits of the action are called the **conjugacy classes** of G. If $x_1, \ldots, x_k$ are representatives from each of the nontrivial conjugacy classes of G and $|\mathcal{O}_{x_1}| = n_1, \ldots, |\mathcal{O}_{x_k}| = n_k$, then

$$|G| = |Z(G)| + n_1 + \cdots + n_k.$$

The stabilizer subgroups of each of the x_i's, $C(x_i) = \{g \in G : gx_i = x_ig\}$, are called the **centralizer subgroups** of the x_i's. From Theorem 14.11, we obtain the **class equation**:

$$|G| = |Z(G)| + [G : C(x_1)] + \cdots + [G : C(x_k)].$$

One of the consequences of the class equation is that the order of each conjugacy class must divide the order of G.

Example 14.12. It is easy to check that the conjugacy classes in S_3 are the following:

$$\{(1)\}, \quad \{(123), (132)\}, \quad \{(12), (13), (23)\}.$$

The class equation is $6 = 1 + 2 + 3$.

Example 14.13. The center of D_4 is $\{(1), (13)(24)\}$, and the conjugacy classes are

$$\{(13), (24)\}, \quad \{(1432), (1234)\}, \quad \{(12)(34), (14)(23)\}.$$

Thus, the class equation for D_4 is $8 = 2 + 2 + 2 + 2$.

Example 14.14. For S_n it takes a bit of work to find the conjugacy classes. We begin with cycles. Suppose that $\sigma = (a_1, \ldots, a_k)$ is a cycle and let $\tau \in S_n$. By Theorem 6.16,

$$\tau\sigma\tau^{-1} = (\tau(a_1), \ldots, \tau(a_k)).$$

Consequently, any two cycles of the same length are conjugate. Now let $\sigma = \sigma_1\sigma_2 \cdots \sigma_r$ be a cycle decomposition, where the length of each cycle σ_i is r_i. Then σ is conjugate to every other $\tau \in S_n$ whose cycle decomposition has the same lengths.

The number of conjugate classes in S_n is the number of ways in which n can be partitioned into sums of positive integers. In the case of S_3 for example, we can partition the integer 3 into the following three sums:

$$3 = 1 + 1 + 1$$
$$3 = 1 + 2$$

$$3 = 3;$$

therefore, there are three conjugacy classes. The problem of finding the number of such partitions for any positive integer n is what computer scientists call **NP-complete**. This effectively means that the problem cannot be solved for a large n because the computations would be too time-consuming for even the largest computer.

Theorem 14.15. *Let G be a group of order p^n where p is prime. Then G has a nontrivial center.*

PROOF.   We apply the class equation

$$|G| = |Z(G)| + n_1 + \cdots + n_k.$$

Since each $n_i > 1$ and $n_i \mid |G|$, it follows that p must divide each n_i. Also, $p \mid |G|$; hence, p must divide $|Z(G)|$. Since the identity is always in the center of G, $|Z(G)| \geq 1$. Therefore, $|Z(G)| \geq p$, and there exists some $g \in Z(G)$ such that $g \neq 1$. $\square$

Corollary 14.16. *Let G be a group of order p^2 where p is prime. Then G is abelian.*

PROOF.   By Theorem 14.15, $|Z(G)| = p$ or p^2. If $|Z(G)| = p^2$, then we are done. Suppose that $|Z(G)| = p$. Then $Z(G)$ and $G/Z(G)$ both have order p and must both be cyclic groups. Choosing a generator $aZ(G)$ for $G/Z(G)$, we can write any element $gZ(G)$ in the quotient group as $a^m Z(G)$ for some integer m; hence, $g = a^m x$ for some x in the center of G. Similarly, if $hZ(G) \in G/Z(G)$, there exists a y in $Z(G)$ such that $h = a^n y$ for some integer n. Since x and y are in the center of G, they commute with all other elements of G; therefore,

$$gh = a^m x a^n y = a^{m+n} xy = a^n y a^m x = hg,$$

and G must be abelian. $\square$

14.3   Burnside's Counting Theorem

Suppose that we wish to color the vertices of a square with two different colors, say black and white. We might suspect that there would be $2^4 = 16$ different colorings. However, some of these colorings are equivalent. If we color the first vertex black and the remaining vertices white, it is the same as coloring the second vertex black and the remaining ones white since we could obtain the second coloring simply by rotating the square $90°$ (Figure 14.17).
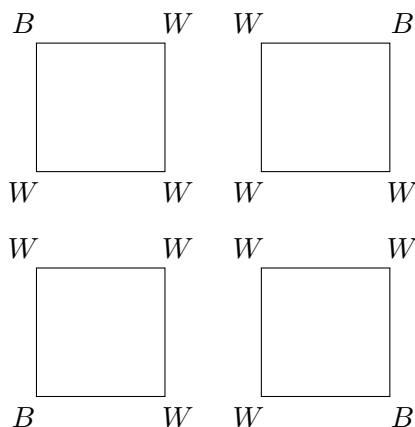

Figure 14.17: Equivalent colorings of square

Burnside's Counting Theorem offers a method of computing the number of distinguishable ways in which something can be done. In addition to its geometric applications, the theorem has interesting applications to areas in switching theory and chemistry. The proof of Burnside's Counting Theorem depends on the following lemma.

Lemma 14.18. *Let X be a G-set and suppose that $x \sim y$. Then G_x is isomorphic to G_y. In particular, $|G_x| = |G_y|$.*

PROOF. Let G act on X by $(g, x) \mapsto g \cdot x$. Since $x \sim y$, there exists a $g \in G$ such that $g \cdot x = y$. Let $a \in G_x$. Since

$$gag^{-1} \cdot y = ga \cdot g^{-1}y = ga \cdot x = g \cdot x = y,$$

we can define a map $\phi : G_x \to G_y$ by $\phi(a) = gag^{-1}$. The map ϕ is a homomorphism since

$$\phi(ab) = gabg^{-1} = gag^{-1}gbg^{-1} = \phi(a)\phi(b).$$

Suppose that $\phi(a) = \phi(b)$. Then $gag^{-1} = gbg^{-1}$ or $a = b$; hence, the map is injective. To show that ϕ is onto, let b be in G_y; then $g^{-1}bg$ is in G_x since

$$g^{-1}bg \cdot x = g^{-1}b \cdot gx = g^{-1}b \cdot y = g^{-1} \cdot y = x;$$

and $\phi(g^{-1}bg) = b$. $\qquad\square$

Theorem 14.19 Burnside. *Let G be a finite group acting on a set X and let k denote the number of orbits of X. Then*

$$k = \frac{1}{|G|} \sum_{g \in G} |X_g|.$$

PROOF. We look at all the fixed points x of all the elements in $g \in G$; that is, we look at all g's and all x's such that $gx = x$. If viewed in terms of fixed point sets, the number of all g's fixing x's is

$$\sum_{g \in G} |X_g|.$$

However, if viewed in terms of the stabilizer subgroups, this number is

$$\sum_{x \in X} |G_x|;$$

hence, $\sum_{g \in G} |X_g| = \sum_{x \in X} |G_x|$. By Lemma 14.18,

$$\sum_{y \in \mathcal{O}_x} |G_y| = |\mathcal{O}_x| \cdot |G_x|.$$

By Theorem 14.11 and Lagrange's Theorem, this expression is equal to $|G|$. Summing over all of the k distinct orbits, we conclude that

$$\sum_{g \in G} |X_g| = \sum_{x \in X} |G_x| = k \cdot |G|. \qquad\square$$

Example 14.20. Let $X = \{1, 2, 3, 4, 5\}$ and suppose that G is the permutation group $G = \{(1), (13), (13)(25), (25)\}$. The orbits of X are $\{1, 3\}$, $\{2, 5\}$, and $\{4\}$. The fixed point sets are

$$X_{(1)} = X$$

$$X_{(13)} = \{2, 4, 5\}$$
$$X_{(13)(25)} = \{4\}$$
$$X_{(25)} = \{1, 3, 4\}.$$

Burnside's Theorem says that

$$k = \frac{1}{|G|} \sum_{g \in G} |X_g| = \frac{1}{4}(5 + 3 + 1 + 3) = 3.$$

A Geometric Example

Before we apply Burnside's Theorem to switching-theory problems, let us examine the number of ways in which the vertices of a square can be colored black or white. Notice that we can sometimes obtain equivalent colorings by simply applying a rigid motion to the square. For instance, as we have pointed out, if we color one of the vertices black and the remaining three white, it does not matter which vertex was colored black since a rotation will give an equivalent coloring.

The symmetry group of a square, D_4, is given by the following permutations:

(1)	(13)	(24)	(1432)
(1234)	(12)(34)	(14)(23)	(13)(24)

The group G acts on the set of vertices $\{1, 2, 3, 4\}$ in the usual manner. We can describe the different colorings by mappings from X into $Y = \{B, W\}$ where B and W represent the colors black and white, respectively. Each map $f : X \to Y$ describes a way to color the corners of the square. Every $\sigma \in D_4$ induces a permutation $\widetilde{\sigma}$ of the possible colorings given by $\widetilde{\sigma}(f) = f \circ \sigma$ for $f : X \to Y$. For example, suppose that f is defined by

$$f(1) = B$$
$$f(2) = W$$
$$f(3) = W$$
$$f(4) = W$$

and $\sigma = (12)(34)$. Then $\widetilde{\sigma}(f) = f \circ \sigma$ sends vertex 2 to B and the remaining vertices to W. The set of all such $\widetilde{\sigma}$ is a permutation group $\widetilde{G}$ on the set of possible colorings. Let $\widetilde{X}$ denote the set of all possible colorings; that is, $\widetilde{X}$ is the set of all possible maps from X to Y. Now we must compute the number of $\widetilde{G}$-equivalence classes.

1. $\widetilde{X}_{(1)} = \widetilde{X}$ since the identity fixes every possible coloring. $|\widetilde{X}| = 2^4 = 16$.

2. $\widetilde{X}_{(1234)}$ consists of all $f \in \widetilde{X}$ such that f is unchanged by the permutation (1234). In this case $f(1) = f(2) = f(3) = f(4)$, so that all values of f must be the same; that is, either $f(x) = B$ or $f(x) = W$ for every vertex x of the square. So $|\widetilde{X}_{(1234)}| = 2$.

3. $|\widetilde{X}_{(1432)}| = 2$.

4. For $\widetilde{X}_{(13)(24)}$, $f(1) = f(3)$ and $f(2) = f(4)$. Thus, $|\widetilde{X}_{(13)(24)}| = 2^2 = 4$.

5. $|\widetilde{X}_{(12)(34)}| = 4$.

6. $|\widetilde{X}_{(14)(23)}| = 4$.

7. For $\widetilde{X}_{(13)}$, $f(1) = f(3)$ and the other corners can be of any color; hence, $|\widetilde{X}_{(13)}| = 2^3 = 8$.

8. $|\widetilde{X}_{(24)}| = 8$.

By Burnside's Theorem, we can conclude that there are exactly

$$\frac{1}{8}(2^4 + 2^1 + 2^2 + 2^1 + 2^2 + 2^2 + 2^3 + 2^3) = 6$$

ways to color the vertices of the square.

Proposition 14.21. *Let G be a permutation group of X and $\widetilde{X}$ the set of functions from X to Y. Then there exists a permutation group $\widetilde{G}$ acting on $\widetilde{X}$, where $\widetilde{\sigma} \in \widetilde{G}$ is defined by $\widetilde{\sigma}(f) = f \circ \sigma$ for $\sigma \in G$ and $f \in \widetilde{X}$. Furthermore, if n is the number of cycles in the cycle decomposition of σ, then $|\widetilde{X}_\sigma| = |Y|^n$.*

PROOF.   Let $\sigma \in G$ and $f \in \widetilde{X}$. Clearly, $f \circ \sigma$ is also in $\widetilde{X}$. Suppose that g is another function from X to Y such that $\widetilde{\sigma}(f) = \widetilde{\sigma}(g)$. Then for each $x \in X$,

$$f(\sigma(x)) = \widetilde{\sigma}(f)(x) = \widetilde{\sigma}(g)(x) = g(\sigma(x)).$$

Since σ is a permutation of X, every element x' in X is the image of some x in X under σ; hence, f and g agree on all elements of X. Therefore, $f = g$ and $\widetilde{\sigma}$ is injective. The map $\sigma \mapsto \widetilde{\sigma}$ is onto, since the two sets are the same size.

Suppose that σ is a permutation of X with cycle decomposition $\sigma = \sigma_1 \sigma_2 \cdots \sigma_n$. Any f in $\widetilde{X}_\sigma$ must have the same value on each cycle of σ. Since there are n cycles and $|Y|$ possible values for each cycle, $|\widetilde{X}_\sigma| = |Y|^n$. □

Example 14.22. Let $X = \{1, 2, \ldots, 7\}$ and suppose that $Y = \{A, B, C\}$. If g is the permutation of X given by $(13)(245) = (13)(245)(6)(7)$, then $n = 4$. Any $f \in \widetilde{X}_g$ must have the same value on each cycle in g. There are $|Y| = 3$ such choices for any value, so $|\widetilde{X}_g| = 3^4 = 81$.

Example 14.23. Suppose that we wish to color the vertices of a square using four different colors. By Proposition 14.21, we can immediately decide that there are

$$\frac{1}{8}(4^4 + 4^1 + 4^2 + 4^1 + 4^2 + 4^2 + 4^3 + 4^3) = 55$$

possible ways.

Switching Functions

In switching theory we are concerned with the design of electronic circuits with binary inputs and outputs. The simplest of these circuits is a switching function that has n inputs and a single output (Figure 14.24). Large electronic circuits can often be constructed by combining smaller modules of this kind. The inherent problem here is that even for a simple circuit a large number of different switching functions can be constructed. With only four inputs and a single output, we can construct 65,536 different switching functions. However, we can often replace one switching function with another merely by permuting the input leads to the circuit (Figure 14.25).
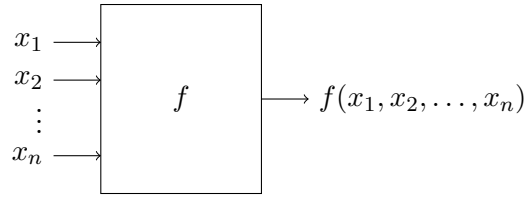
Figure 14.24: A switching function of n variables

We define a ***switching*** or ***Boolean function*** of n variables to be a function from $\mathbb{Z}_2^n$ to $\mathbb{Z}_2$. Since any switching function can have two possible values for each binary n-tuple and there are 2^n binary n-tuples, 2^{2^n} switching functions are possible for n variables. In general, allowing permutations of the inputs greatly reduces the number of different kinds of modules that are needed to build a large circuit.
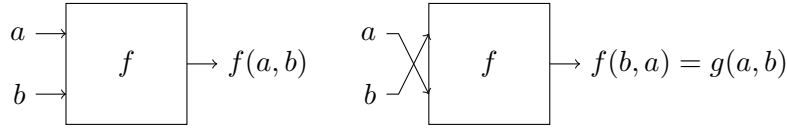

Figure 14.25: A switching function of two variables

The possible switching functions with two input variables a and b are listed in Table 14.26. Two switching functions f and g are equivalent if g can be obtained from f by a permutation of the input variables. For example, $g(a, b, c) = f(b, c, a)$. In this case $g \sim f$ via the permutation (acb). In the case of switching functions of two variables, the permutation (ab) reduces 16 possible switching functions to 12 equivalent functions since

$$f_2 \sim f_4$$
$$f_3 \sim f_5$$
$$f_{10} \sim f_{12}$$
$$f_{11} \sim f_{13}.$$

Inputs		Outputs							
		f_0	f_1	f_2	f_3	f_4	f_5	f_6	f_7
0	0	0	0	0	0	0	0	0	0
0	1	0	0	0	0	1	1	1	1
1	0	0	0	1	1	0	0	1	1
1	1	0	1	0	1	0	1	0	1
Inputs		Outputs							
		f_8	f_9	f_{10}	f_{11}	f_{12}	f_{13}	f_{14}	f_{15}
0	0	1	1	1	1	1	1	1	1
0	1	0	0	0	0	1	1	1	1
1	0	0	0	1	1	0	0	1	1
1	1	0	1	0	1	0	1	0	1

Table 14.26: Switching functions in two variables

For three input variables there are $2^{2^3} = 256$ possible switching functions; in the case of four variables there are $2^{2^4} = 65{,}536$. The number of equivalence classes is too large to

reasonably calculate directly. It is necessary to employ Burnside's Theorem.

Consider a switching function with three possible inputs, a, b, and c. As we have mentioned, two switching functions f and g are equivalent if a permutation of the input variables of f gives g. It is important to notice that a permutation of the switching functions is not simply a permutation of the input values $\{a, b, c\}$. A switching function is a set of output values for the inputs a, b, and c, so when we consider equivalent switching functions, we are permuting 2^3 possible outputs, not just three input values. For example, each binary triple (a, b, c) has a specific output associated with it. The permutation (acb) changes outputs as follows:

$$(0, 0, 0) \mapsto (0, 0, 0)$$
$$(0, 0, 1) \mapsto (0, 1, 0)$$
$$(0, 1, 0) \mapsto (1, 0, 0)$$
$$\vdots$$
$$(1, 1, 0) \mapsto (1, 0, 1)$$
$$(1, 1, 1) \mapsto (1, 1, 1).$$

Let X be the set of output values for a switching function in n variables. Then $|X| = 2^n$. We can enumerate these values as follows:

$$(0, \ldots, 0, 1) \mapsto 0$$
$$(0, \ldots, 1, 0) \mapsto 1$$
$$(0, \ldots, 1, 1) \mapsto 2$$
$$\vdots$$
$$(1, \ldots, 1, 1) \mapsto 2^n - 1.$$

Now let us consider a circuit with four input variables and a single output. Suppose that we can permute the leads of any circuit according to the following permutation group:

$$(a), \quad (ac), \quad (bd), \quad (adcb),$$
$$(abcd), \quad (ab)(cd), \quad (ad)(bc), \quad (ac)(bd).$$

The permutations of the four possible input variables induce the permutations of the output values in Table 14.27.

Hence, there are

$$\frac{1}{8}(2^{16} + 2 \cdot 2^{12} + 2 \cdot 2^6 + 3 \cdot 2^{10}) = 9616$$

possible switching functions of four variables under this group of permutations. This number will be even smaller if we consider the full symmetric group on four letters.

Group Permutation	Switching Function Permutation	Number of Cycles
(a)	(0)	16
(ac)	$(2,8)(3,9)(6,12)(7,13)$	12
(bd)	$(1,4)(3,6)(9,12)(11,14)$	12
$(adcb)$	$(1,2,4,8)(3,6.12,9)(5,10)(7,14,13,11)$	6
$(abcd)$	$(1,8,4,2)(3,9,12,6)(5,10)(7,11,13,14)$	6
$(ab)(cd)$	$(1,2)(4,8)(5,10)(6,9)(7,11)(13,14)$	10
$(ad)(bc)$	$(1,8)(2,4)(3,12)(5,10)(7,14)(11,13)$	10
$(ac)(bd)$	$(1,4)(2,8)(3,12)(6,9)(7,13)(11,14)$	10

Table 14.27: Permutations of switching functions in four variables

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Historical Note

William Burnside was born in London in 1852. He attended Cambridge University from 1871 to 1875 and won the Smith's Prize in his last year. After his graduation he lectured at Cambridge. He was made a member of the Royal Society in 1893. Burnside wrote approximately 150 papers on topics in applied mathematics, differential geometry, and probability, but his most famous contributions were in group theory. Several of Burnside's conjectures have stimulated research to this day. One such conjecture was that every group of odd order is solvable; that is, for a group G of odd order, there exists a sequence of subgroups

$$G = H_n \supset H_{n-1} \supset \cdots \supset H_1 \supset H_0 = \{e\}$$

such that H_i is normal in H_{i+1} and H_{i+1}/H_i is abelian. This conjecture was finally proven by W. Feit and J. Thompson in 1963. Burnside's *The Theory of Groups of Finite Order*, published in 1897, was one of the first books to treat groups in a modern context as opposed to permutation groups. The second edition, published in 1911, is still a classic.

14.4 Exercises

1.   Examples 14.1–14.5 in the first section each describe an action of a group G on a set X, which will give rise to the equivalence relation defined by G-equivalence. For each example, compute the equivalence classes of the equivalence relation, the G-**equivalence classes**.

2.   Compute all X_g and all G_x for each of the following permutation groups.

  (a) $X = \{1,2,3\}$, $G = S_3 = \{(1),(12),(13),(23),(123),(132)\}$

  (b) $X = \{1,2,3,4,5,6\}$, $G = \{(1),(12),(345),(354),(12)(345),(12)(354)\}$

3.   Compute the G-equivalence classes of X for each of the G-sets in Exercise 14.4.2. For each $x \in X$ verify that $|G| = |\mathcal{O}_x| \cdot |G_x|$.

4.   Let G be the additive group of real numbers. Let the action of $\theta \in G$ on the real plane $\mathbb{R}^2$ be given by rotating the plane counterclockwise about the origin through θ radians. Let P be a point on the plane other than the origin.

  (a) Show that $\mathbb{R}^2$ is a G-set.

  (b) Describe geometrically the orbit containing P.

(c) Find the group G_P.

5. Let $G = A_4$ and suppose that G acts on itself by conjugation; that is, $(g, h) \mapsto ghg^{-1}$.

(a) Determine the conjugacy classes (orbits) of each element of G.

(b) Determine all of the isotropy subgroups for each element of G.

6. Find the conjugacy classes and the class equation for each of the following groups.

(a) S_4 (b) D_5 (c) $\mathbb{Z}_9$ (d) Q_8

7. Write the class equation for S_5 and for A_5.

8. If a square remains fixed in the plane, how many different ways can the corners of the square be colored if three colors are used?

9. How many ways can the vertices of an equilateral triangle be colored using three different colors?

10. Find the number of ways a six-sided die can be constructed if each side is marked differently with $1, \ldots, 6$ dots.

11. Up to a rotation, how many ways can the faces of a cube be colored with three different colors?

12. Consider 12 straight wires of equal lengths with their ends soldered together to form the edges of a cube. Either silver or copper wire can be used for each edge. How many different ways can the cube be constructed?

13. Suppose that we color each of the eight corners of a cube. Using three different colors, how many ways can the corners be colored up to a rotation of the cube?

14. Each of the faces of a regular tetrahedron can be painted either red or white. Up to a rotation, how many different ways can the tetrahedron be painted?

15. Suppose that the vertices of a regular hexagon are to be colored either red or white. How many ways can this be done up to a symmetry of the hexagon?

16. A molecule of benzene is made up of six carbon atoms and six hydrogen atoms, linked together in a hexagonal shape as in Figure 14.28.

(a) How many different compounds can be formed by replacing one or more of the hydrogen atoms with a chlorine atom?

(b) Find the number of different chemical compounds that can be formed by replacing three of the six hydrogen atoms in a benzene ring with a CH_3 radical.
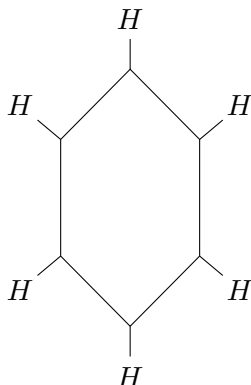

Figure 14.28: A benzene ring

17.   How many equivalence classes of switching functions are there if the input variables x_1, x_2, and x_3 can be permuted by any permutation in S_3? What if the input variables x_1, x_2, x_3, and x_4 can be permuted by any permutation in S_4?

18.   How many equivalence classes of switching functions are there if the input variables x_1, x_2, x_3, and x_4 can be permuted by any permutation in the subgroup of S_4 generated by the permutation $(x_1 x_2 x_3 x_4)$?

19.   A striped necktie has 12 bands of color. Each band can be colored by one of four possible colors. How many possible different-colored neckties are there?

20.   A group acts ***faithfully*** on a G-set X if the identity is the only element of G that leaves every element of X fixed. Show that G acts faithfully on X if and only if no two distinct elements of G have the same action on each element of X.

21.   Let p be prime. Show that the number of different abelian groups of order p^n (up to isomorphism) is the same as the number of conjugacy classes in S_n.

22.   Let $a \in G$. Show that for any $g \in G$, $gC(a)g^{-1} = C(gag^{-1})$.

23.   Let $|G| = p^n$ be a nonabelian group for p prime. Prove that $|Z(G)| < p^{n-1}$.

24.   Let G be a group with order p^n where p is prime and X a finite G-set. If $X_G = \{x \in X : gx = x \text{ for all } g \in G\}$ is the set of elements in X fixed by the group action, then prove that $|X| \equiv |X_G| \pmod{p}$.

25.   If G is a group of order p^n, where p is prime and $n \geq 2$, show that G must have a proper subgroup of order p. If $n \geq 3$, is it true that G will have a proper subgroup of order p^2?

14.5   Programming Exercise

1.   Write a program to compute the number of conjugacy classes in S_n. What is the largest n for which your program will work?

References and Suggested Reading

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14.6 Sage

Groups can be realized in many ways, such as as sets of permutations, as sets of matrices, or as sets of abstract symbols related by certain rules ("presentations") and in myriad other ways. We have concentrated on permutation groups because of their concrete feel, with elements written as functions, and because of their thorough implementation in Sage. Group actions are of great interest when the set they act on is the group itself, and group actions will figure prominently in the proofs of the main results of the next chapter. However, any time we have a group action on a set, we can view that group as a permutation group on the elements of the set. So permutation groups are an area of group theory of independent interest, with its own definitions and theorems.

We will describe Sage's commands applicable when a group action arises naturally via conjugation, and then move into the more general situation in a more general application.

Conjugation as a Group Action

We might think we need to be careful how Sage defines conjugation (gxg^{-1} versus $g^{-1}xg$) and the difference between Sage and the text on the order of products. However, if you look at the definition of the center and centralizer subgroups you can see that any difference in ordering is irrelevant. Here are the group action commands for the particular action that is conjugation of the elements of the group.

Sage has a permutation group method `.center()` which returns the subgroup of fixed points. The permutation group method, `.centralizer(g)`, returns a subgroup that is the stabilizer of the group element `g`. Finally, the orbits are given by conjugacy classes, but Sage will not flood you with the full conjugacy classes and instead gives back a list of one element per conjugacy class, the representatives, via the permutation group method `.conjugacy_classes_representatives()`. You can manually reconstruct a conjugacy class from a representative, as we do in the example below.

Here is an example of the above commands in action. Notice that an abelian group would be a bad choice for this example.

```
D = DihedralGroup(8)
C = D.center(); C
```

```
Subgroup of (Dihedral group of order 16 as a permutation group)
generated by [(1,5)(2,6)(3,7)(4,8)]
```

```
C.list()
```

```
[(), (1,5)(2,6)(3,7)(4,8)]
```

```
a = D("(1,2)(3,8)(4,7)(5,6)")
C1 = D.centralizer(a); C1.list()
```

```
[(), (1,2)(3,8)(4,7)(5,6), (1,5)(2,6)(3,7)(4,8), (1,6)(2,5)(3,4)(7,8)]
```

```
b = D("(1,2,3,4,5,6,7,8)")
C2 = D.centralizer(b); C2.order()
```

```
8
```

```
CCR = D.conjugacy_classes_representatives(); CCR
```

```
[(), (2,8)(3,7)(4,6), (1,2)(3,8)(4,7)(5,6), (1,2,3,4,5,6,7,8),
 (1,3,5,7)(2,4,6,8), (1,4,7,2,5,8,3,6), (1,5)(2,6)(3,7)(4,8)]
```

```
r = CCR[2]; r
```

```
(1,2)(3,8)(4,7)(5,6)
```

```
conj = []
x = [conj.append(g^-1*r*g) for g in D if not g^-1*r*g in conj]
conj
```

```
[(1,2)(3,8)(4,7)(5,6), (1,4)(2,3)(5,8)(6,7),
 (1,6)(2,5)(3,4)(7,8), (1,8)(2,7)(3,6)(4,5)]
```

Notice that in the one conjugacy class constructed all the elements have the same cycle structure, which is no accident. Notice too that rep and a are the same element, and the product of the order of the centralizer (4) and the size of the conjugacy class (4) equals the order of the group (16), which is a variant of the conclusion of Theorem 14.11.

Verify that the following is a demonstration of the class equation in the special case when the action is conjugation, but would be valid for any group, rather than just D.

```
sizes = [D.order()/D.centralizer(g).order()
            for g in D.conjugacy_classes_representatives()]
sizes
```

```
[1, 4, 4, 2, 2, 2, 1]
```

```
D.order() == sum(sizes)
```

```
True
```

Graph Automorphisms

As mentioned, group actions can be even more interesting when the set they act on is different from the group itself. One class of examples is the group of symmetries of a geometric solid, where the objects in the set are the vertices of the object, or perhaps some other aspect such as edges, faces or diagonals. In this case, the group is all those permutations that move the solid but leave it filling the same space before the motion ("rigid motions").

In this section we will examine something very similar. A **graph** is a mathematical object, consisting of vertices and edges, but the only structure is whether or not any given pair of vertices are joined by an edge or not. The group consists of permutations of vertices that preserve the structure, that is, permutations of vertices that take edges to edges and non-edges to non-edges. It is very similar to a symmetry group, but there is no notion of any geometric relationships being preserved.

Here is an example. You will need to run the first compute cell to define the graph and get a nice graphic representation.

```
Q = graphs.CubeGraph(3)
Q.plot(layout='spring')
```

```
A = Q.automorphism_group()
A.order()
```

Your plot should look like the vertices and edges of a cube, but may not quite look regular, which is fine, since the geometry is not relevant. Vertices are labeled with strings of three binary digits, 0 or 1, and any two vertices are connected by an edge if their strings differ in exactly one location. We might expect the group of symmetries to have order 24, rather than order 48, given its resemblance to a cube (in appearance and in name). However, when not restricted to rigid motions, we have new permutations that preserve edges. One in particular is to interchange two "opposite faces." Locate two 4-cycles opposite of each other, listed in the same order: $000, 010, 110, 100$ and $001, 011, 111, 101$. Notice that each cycle looks very similar, but all the vertices of the first end in a zero and the second cycle has vertices ending in a one.

  We can create explicitly the permutation that interchanges these two opposite faces, using a text version of the permutation in cycle notation.

```
a = A("('000','001')('010','011')('110','111')('100','101')")
a in A
```

True

We can use this group to illustrate the relevant Sage commands for group actions.

```
A.orbits()
```

```
[['000', '001', '010', '100', '011', '101', '110', '111']]
```

```
len(A.orbits()[0])
```

8

So this action has only one (big) orbit. This implies that every vertex is "like" any other. When a permutation group behaves this way, we say the group is ***transitive***.

```
A.is_transitive()
```

True

If every vertex is "the same" we can compute the stabilizer of any vertex, since they will all be isomorphic. Because vertex 000 is the simplest in some sense, we compute its stabilizer.

```
S = A.stabilizer('000')
S.list()
```

```
[(),
 ('001','100','010')('011','101','110'),
 ('010','100')('011','101'),
 ('001','010','100')('011','110','101'),
 ('001','100')('011','110'),
 ('001','010')('101','110')]
```

That S has 6 elements is no surprise, since the group has order 48 and the size of the lone orbit is 8. But we can go one step further. The three vertices of the graph attached directly to 000 are 100, 010, 001. Any automorphism of the graph that fixes 000 must then permute the three adjacent vertices. There are $3! = 6$ possible ways to do this, and you can check that each appears in one of the six elements of the stabilizer. So we can understand a transitive group by considering the smaller stabilizer, and in this case we can see that each element of the stabilizer is determined by how it permutes the neighbors of the stabilized vertex.

  Transitive groups are both unusual and important. To contrast, here is a graph automorphism group that is far from transitive (without being trivial). A path is a graph that has all of its vertices in a line. Run the first compute cell to see a path on 11 vertices.

```
P = graphs.PathGraph(11)
P.plot()
```

```
A = P.automorphism_group()
A.list()
```

```
[(), (0,10)(1,9)(2,8)(3,7)(4,6)]
```

The automorphism group is the trivial identity automorphism (always) and an order 2 permutation that "flips" the path end-to-end. The group is far from transitive and there are many orbits.

```
A.is_transitive()
```

```
False
```

```
A.orbits()
```

```
[[0, 10], [1, 9], [2, 8], [3, 7], [4, 6], [5]]
```

Most of the stabilizers are trivial, with one exception. As subgroups of a group of order 2, there really are not too many options.

```
A.stabilizer(2).list()
```

```
[()]
```

```
A.stabilizer(5).list()
```

```
[(), (0,10)(1,9)(2,8)(3,7)(4,6)]
```

How would this final example have been different if we had used a path on 10 vertices?

NOTE: There was once a small bug with stabilizers being created as subgroups of symmetric groups on fewer symbols than the correct number. This is fixed in Sage 4.8 and newer. Note the correct output below, and you can check your installation by running these commands. If you do not see the singleton [4] in your output, you should definitely update your copy of Sage.

```
G = SymmetricGroup(4)
S = G.stabilizer(4)
S.orbits()
```

```
[[1, 2, 3], [4]]
```

14.7 Sage Exercises

1.   Construct the Higman-Sims graph with the command `graphs.HigmanSimsGraph()`. Then construct the automorphism group and determine the order of the one interesting normal subgroup of this group. You can try plotting the graph, but the graphic is unlikely to be very informative.

2.   This exercise asks you to verify the class equation outside of the usual situation where the group action is conjugation. Consider the example of the automorphism group of the path on 11 vertices. First construct the list of orbits. From each orbit, grab the first element of the orbit as a representative. Compute the size of the orbit as the index of the stabilizer

of the representative in the group via Theorem 14.11. (Yes, you could just compute the size of the full orbit, but the idea of the exercise is to use more group-theoretic results.) Then sum these orbit-sizes, which should equal the size of the whole vertex set since the orbits form a partition.

3.   Construct a simple graph (no loops or multiple edges), with at least two vertices and at least one edge, whose automorphism group is trivial. You might start experimenting by drawing pictures on paper before constructing the graph. A command like the following will let you construct a graph from edges. The graph below looks like a triangle or 3-cycle.

```
G = Graph([(1,2), (2,3), (3,1)])
G.plot()
```

4.   For the following two pairs of groups, compute the list of conjugacy class representatives for each group in the pair. For each part, compare and contrast the results for the two groups in the pair, with thoughtful and insightful comments.

(a) The full symmetric group on 5 symbols, S_5, and the alternating group on 5 symbols, A_5.

(b) The dihedral groups that are symmetries of a 7-gon and an 8-gon, D_7 and D_8.

5.   Use the command `graphs.CubeGraph(4)` to build the four-dimensional cube graph, Q_4. Using a plain `.plot()` command (without a spring layout) should create a nice plot. Construct the automorphism group of the graph, which will provide a group action on the vertex set.

(a) Construct the orbits of this action, and comment.

(b) Construct a stabilizer of a single vertex (which is a subgroup of the full automorphism group) and then consider the action of *this* group on the vertex set. Construct the orbits of this new action, and comment carefully and fully on your observations, especially in terms of the vertices of the graph.

6.   Build the graph given by the commands below. The result should be a symmetric-looking graph with an automorphism group of order 16.

```
G = graphs.CycleGraph(8)
G.add_edges([(0,2),(1,3),(4,6),(5,7)])
G.plot()
```

Repeat the two parts of the previous exercise, but realize that in the second part there are now two different stabilizers to create, so build both and compare the differences in the stabilizers and their orbits. Creating a second plot with `G.plot(layout='planar')` might provide extra insight.

The Sylow Theorems

We already know that the converse of Lagrange's Theorem is false. If G is a group of order m and n divides m, then G does not necessarily possess a subgroup of order n. For example, A_4 has order 12 but does not possess a subgroup of order 6. However, the Sylow Theorems do provide a partial converse for Lagrange's Theorem—in certain cases they guarantee us subgroups of specific orders. These theorems yield a powerful set of tools for the classification of all finite nonabelian groups.

15.1   The Sylow Theorems

We will use what we have learned about group actions to prove the Sylow Theorems. Recall for a moment what it means for G to act on itself by conjugation and how conjugacy classes are distributed in the group according to the class equation, discussed in Chapter 14. A group G acts on itself by conjugation via the map $(g, x) \mapsto g x g^{-1}$. Let $x_1, \ldots, x_k$ be representatives from each of the distinct conjugacy classes of G that consist of more than one element. Then the class equation can be written as

$$|G| = |Z(G)| + [G : C(x_1)] + \cdots + [G : C(x_k)],$$

where $Z(G) = \{g \in G : gx = xg \text{ for all } x \in G\}$ is the center of G and $C(x_i) = \{g \in G : gx_i = x_ig\}$ is the centralizer subgroup of x_i.

We begin our investigation of the Sylow Theorems by examining subgroups of order p, where p is prime. A group G is a *p-group* if every element in G has as its order a power of p, where p is a prime number. A subgroup of a group G is a *p-subgroup* if it is a p-group.

Theorem 15.1 Cauchy. *Let G be a finite group and p a prime such that p divides the order of G. Then G contains a subgroup of order p.*

PROOF.   We will use induction on the order of G. If $|G| = p$, then clearly G itself is the required subgroup. We now assume that every group of order k, where $p \leq k < n$ and p divides k, has an element of order p. Assume that $|G| = n$ and $p \mid n$ and consider the class equation of G:

$$|G| = |Z(G)| + [G : C(x_1)] + \cdots + [G : C(x_k)].$$

We have two cases.

Case 1. Suppose the order of one of the centralizer subgroups, $C(x_i)$, is divisible by p for some i, $i = 1, \ldots, k$. In this case, by our induction hypothesis, we are done. Since $C(x_i)$ is a proper subgroup of G and p divides $|C(x_i)|$, $C(x_i)$ must contain an element of order p. Hence, G must contain an element of order p.

Case 2. Suppose the order of no centralizer subgroup is divisible by p. Then p divides $[G : C(x_i)]$, the order of each conjugacy class in the class equation; hence, p must divide

the center of G, $Z(G)$. Since $Z(G)$ is abelian, it must have a subgroup of order p by the Fundamental Theorem of Finite Abelian Groups. Therefore, the center of G contains an element of order p. $\square$

Corollary 15.2. *Let G be a finite group. Then G is a p-group if and only if $|G| = p^n$.*

Example 15.3. Let us consider the group A_5. We know that $|A_5| = 60 = 2^2 \cdot 3 \cdot 5$. By Cauchy's Theorem, we are guaranteed that A_5 has subgroups of orders 2, 3 and 5. The Sylow Theorems will give us even more information about the possible subgroups of A_5.

We are now ready to state and prove the first of the Sylow Theorems. The proof is very similar to the proof of Cauchy's Theorem.

Theorem 15.4 First Sylow Theorem. *Let G be a finite group and p a prime such that p^r divides $|G|$. Then G contains a subgroup of order p^r.*

PROOF.   We induct on the order of G once again. If $|G| = p$, then we are done. Now suppose that the order of G is n with $n > p$ and that the theorem is true for all groups of order less than n, where p divides n. We shall apply the class equation once again:

$$|G| = |Z(G)| + [G : C(x_1)] + \cdots + [G : C(x_k)].$$

First suppose that p does not divide $[G : C(x_i)]$ for some i. Then $p^r \mid |C(x_i)|$, since p^r divides $|G| = |C(x_i)| \cdot [G : C(x_i)]$. Now we can apply the induction hypothesis to $C(x_i)$.

Hence, we may assume that p divides $[G : C(x_i)]$ for all i. Since p divides $|G|$, the class equation says that p must divide $|Z(G)|$; hence, by Cauchy's Theorem, $Z(G)$ has an element of order p, say g. Let N be the group generated by g. Clearly, N is a normal subgroup of $Z(G)$ since $Z(G)$ is abelian; therefore, N is normal in G since every element in $Z(G)$ commutes with every element in G. Now consider the factor group G/N of order $|G|/p$. By the induction hypothesis, G/N contains a subgroup H of order p^{r-1}. The inverse image of H under the canonical homomorphism $\phi : G \to G/N$ is a subgroup of order p^r in G. $\square$

A **Sylow p-subgroup** P of a group G is a maximal p-subgroup of G. To prove the other two Sylow Theorems, we need to consider conjugate subgroups as opposed to conjugate elements in a group. For a group G, let $\mathcal{S}$ be the collection of all subgroups of G. For any subgroup H, $\mathcal{S}$ is a H-set, where H acts on $\mathcal{S}$ by conjugation. That is, we have an action

$$H \times \mathcal{S} \to \mathcal{S}$$

defined by

$$h \cdot K \mapsto hKh^{-1}$$

for K in $\mathcal{S}$.

The set

$$N(H) = \{g \in G : gHg^{-1} = H\}$$

is a subgroup of G called the the **normalizer** of H in G. Notice that H is a normal subgroup of $N(H)$. In fact, $N(H)$ is the largest subgroup of G in which H is normal.

Lemma 15.5. *Let P be a Sylow p-subgroup of a finite group G and let x have as its order a power of p. If $x^{-1}Px = P$, then $x \in P$.*

PROOF.   Certainly $x \in N(P)$, and the cyclic subgroup, $\langle xP \rangle \subset N(P)/P$, has as its order a power of p. By the Correspondence Theorem there exists a subgroup H of $N(P)$ containing P such that $H/P = \langle xP \rangle$. Since $|H| = |P| \cdot |\langle xP \rangle|$, the order of H must be a power of p. However, P is a Sylow p-subgroup contained in H. Since the order of P is the largest power of p dividing $|G|$, $H = P$. Therefore, H/P is the trivial subgroup and $xP = P$, or $x \in P$. $\square$

Lemma 15.6. *Let H and K be subgroups of G. The number of distinct H-conjugates of K is $[H : N(K) \cap H]$.*

PROOF. We define a bijection between the conjugacy classes of K and the right cosets of $N(K) \cap H$ by $h^{-1}Kh \mapsto (N(K) \cap H)h$. To show that this map is a bijection, let $h_1, h_2 \in H$ and suppose that $(N(K) \cap H)h_1 = (N(K) \cap H)h_2$. Then $h_2 h_1^{-1} \in N(K)$. Therefore, $K = h_2 h_1^{-1} K h_1 h_2^{-1}$ or $h_1^{-1} K h_1 = h_2^{-1} K h_2$, and the map is an injection. It is easy to see that this map is surjective; hence, we have a one-to-one and onto map between the H-conjugates of K and the right cosets of $N(K) \cap H$ in H. $\square$

Theorem 15.7 Second Sylow Theorem. *Let G be a finite group and p a prime dividing $|G|$. Then all Sylow p-subgroups of G are conjugate. That is, if P_1 and P_2 are two Sylow p-subgroups, there exists a $g \in G$ such that $gP_1 g^{-1} = P_2$.*

PROOF. Let P be a Sylow p-subgroup of G and suppose that $|G| = p^r m$ with $|P| = p^r$. Let

$$\mathcal{S} = \{P = P_1, P_2, \ldots, P_k\}$$

consist of the distinct conjugates of P in G. By Lemma 15.6, $k = [G : N(P)]$. Notice that

$$|G| = p^r m = |N(P)| \cdot [G : N(P)] = |N(P)| \cdot k.$$

Since p^r divides $|N(P)|$, p cannot divide k.

Given any other Sylow p-subgroup Q, we must show that $Q \in \mathcal{S}$. Consider the Q-conjugacy classes of each P_i. Clearly, these conjugacy classes partition $\mathcal{S}$. The size of the partition containing P_i is $[Q : N(P_i) \cap Q]$ by Lemma 15.6, and Lagrange's Theorem tells us that $|Q| = [Q : N(P_i) \cap Q]|N(P_i) \cap Q|$. Thus, $[Q : N(P_i) \cap Q]$ must be a divisor of $|Q| = p^r$. Hence, the number of conjugates in every equivalence class of the partition is a power of p. However, since p does not divide k, one of these equivalence classes must contain only a single Sylow p-subgroup, say P_j. In this case, $x^{-1} P_j x = P_j$ for all $x \in Q$. By Lemma 15.5, $P_j = Q$. $\square$

Theorem 15.8 Third Sylow Theorem. *Let G be a finite group and let p be a prime dividing the order of G. Then the number of Sylow p-subgroups is congruent to $1 \pmod{p}$ and divides $|G|$.*

PROOF. Let P be a Sylow p-subgroup acting on the set of Sylow p-subgroups,

$$\mathcal{S} = \{P = P_1, P_2, \ldots, P_k\},$$

by conjugation. From the proof of the Second Sylow Theorem, the only P-conjugate of P is itself and the order of the other P-conjugacy classes is a power of p. Each P-conjugacy class contributes a positive power of p toward $|\mathcal{S}|$ except the equivalence class $\{P\}$. Since $|\mathcal{S}|$ is the sum of positive powers of p and 1, $|\mathcal{S}| \equiv 1 \pmod{p}$.

Now suppose that G acts on $\mathcal{S}$ by conjugation. Since all Sylow p-subgroups are conjugate, there can be only one orbit under this action. For $P \in \mathcal{S}$,

$$|\mathcal{S}| = |\text{orbit of } P| = [G : N(P)]$$

by Lemma 15.6. But $[G : N(P)]$ is a divisor of $|G|$; consequently, the number of Sylow p-subgroups of a finite group must divide the order of the group. $\square$

Peter Ludvig Mejdell Sylow was born in 1832 in Christiania, Norway (now Oslo). After attending Christiania University, Sylow taught high school. In 1862 he obtained a temporary appointment at Christiania University. Even though his appointment was relatively brief, he influenced students such as Sophus Lie (1842–1899). Sylow had a chance at a permanent chair in 1869, but failed to obtain the appointment. In 1872, he published a 10-page paper presenting the theorems that now bear his name. Later Lie and Sylow collaborated on a new edition of Abel's works. In 1898, a chair at Christiania University was finally created for Sylow through the efforts of his student and colleague Lie. Sylow died in 1918.

15.2 Examples and Applications

Example 15.9. Using the Sylow Theorems, we can determine that A_5 has subgroups of orders 2, 3, 4, and 5. The Sylow p-subgroups of A_5 have orders 3, 4, and 5. The Third Sylow Theorem tells us exactly how many Sylow p-subgroups A_5 has. Since the number of Sylow 5-subgroups must divide 60 and also be congruent to 1 (mod 5), there are either one or six Sylow 5-subgroups in A_5. All Sylow 5-subgroups are conjugate. If there were only a single Sylow 5-subgroup, it would be conjugate to itself; that is, it would be a normal subgroup of A_5. Since A_5 has no normal subgroups, this is impossible; hence, we have determined that there are exactly six distinct Sylow 5-subgroups of A_5.

The Sylow Theorems allow us to prove many useful results about finite groups. By using them, we can often conclude a great deal about groups of a particular order if certain hypotheses are satisfied.

Theorem 15.10. *If p and q are distinct primes with $p < q$, then every group G of order pq has a single subgroup of order q and this subgroup is normal in G. Hence, G cannot be simple. Furthermore, if $q \not\equiv 1 \pmod{p}$, then G is cyclic.*

PROOF. We know that G contains a subgroup H of order q. The number of conjugates of H divides pq and is equal to $1 + kq$ for $k = 0, 1, \ldots$. However, $1 + q$ is already too large to divide the order of the group; hence, H can only be conjugate to itself. That is, H must be normal in G.

The group G also has a Sylow p-subgroup, say K. The number of conjugates of K must divide q and be equal to $1 + kp$ for $k = 0, 1, \ldots$. Since q is prime, either $1 + kp = q$ or $1 + kp = 1$. If $1 + kp = 1$, then K is normal in G. In this case, we can easily show that G satisfies the criteria, given in Chapter 9, for the internal direct product of H and K. Since H is isomorphic to $\mathbb{Z}_q$ and K is isomorphic to $\mathbb{Z}_p$, $G \cong \mathbb{Z}_p \times \mathbb{Z}_q \cong \mathbb{Z}_{pq}$ by Theorem 9.21. □

Example 15.11. Every group of order 15 is cyclic. This is true because $15 = 5 \cdot 3$ and $5 \not\equiv 1 \pmod{3}$.

Example 15.12. Let us classify all of the groups of order $99 = 3^2 \cdot 11$ up to isomorphism. First we will show that every group G of order 99 is abelian. By the Third Sylow Theorem, there are $1+3k$ Sylow 3-subgroups, each of order 9, for some $k = 0, 1, 2, \ldots$. Also, $1+3k$ must divide 11; hence, there can only be a single normal Sylow 3-subgroup H in G. Similarly, there are $1+11k$ Sylow 11-subgroups and $1+11k$ must divide 9. Consequently, there is only one Sylow 11-subgroup K in G. By Corollary 14.16, any group of order p^2 is abelian for p prime; hence, H is isomorphic either to $\mathbb{Z}_3 \times \mathbb{Z}_3$ or to $\mathbb{Z}_9$. Since K has order 11, it must be isomorphic to $\mathbb{Z}_{11}$. Therefore, the only possible groups of order 99 are $\mathbb{Z}_3 \times \mathbb{Z}_3 \times \mathbb{Z}_{11}$ or $\mathbb{Z}_9 \times \mathbb{Z}_{11}$ up to isomorphism.

To determine all of the groups of order $5 \cdot 7 \cdot 47 = 1645$, we need the following theorem.

Theorem 15.13. *Let $G' = \langle aba^{-1}b^{-1} : a, b \in G \rangle$ be the subgroup consisting of all finite products of elements of the form $aba^{-1}b^{-1}$ in a group G. Then G' is a normal subgroup of G and G/G' is abelian.*

The subgroup G' of G is called the **commutator subgroup** of G. We leave the proof of this theorem as an exercise (Exercise 10.3.14 in Chapter 10).

Example 15.14. We will now show that every group of order $5 \cdot 7 \cdot 47 = 1645$ is abelian, and cyclic by Corollary 9.21. By the Third Sylow Theorem, G has only one subgroup H_1 of order 47. So G/H_1 has order 35 and must be abelian by Theorem 15.10. Hence, the commutator subgroup of G is contained in H which tells us that $|G'|$ is either 1 or 47. If $|G'| = 1$, we are done. Suppose that $|G'| = 47$. The Third Sylow Theorem tells us that G has only one subgroup of order 5 and one subgroup of order 7. So there exist normal subgroups H_2 and H_3 in G, where $|H_2| = 5$ and $|H_3| = 7$. In either case the quotient group is abelian; hence, G' must be a subgroup of H_i, $i = 1, 2$. Therefore, the order of G' is 1, 5, or 7. However, we already have determined that $|G'| = 1$ or 47. So the commutator subgroup of G is trivial, and consequently G is abelian.

Finite Simple Groups

Given a finite group, one can ask whether or not that group has any normal subgroups. Recall that a simple group is one with no proper nontrivial normal subgroups. As in the case of A_5, proving a group to be simple can be a very difficult task; however, the Sylow Theorems are useful tools for proving that a group is not simple. Usually, some sort of counting argument is involved.

Example 15.15. Let us show that no group G of order 20 can be simple. By the Third Sylow Theorem, G contains one or more Sylow 5-subgroups. The number of such subgroups is congruent to 1 (mod 5) and must also divide 20. The only possible such number is 1. Since there is only a single Sylow 5-subgroup and all Sylow 5-subgroups are conjugate, this subgroup must be normal.

Example 15.16. Let G be a finite group of order p^n, $n > 1$ and p prime. By Theorem 14.15, G has a nontrivial center. Since the center of any group G is a normal subgroup, G cannot be a simple group. Therefore, groups of orders 4, 8, 9, 16, 25, 27, 32, 49, 64, and 81 are not simple. In fact, the groups of order 4, 9, 25, and 49 are abelian by Corollary 14.16.

Example 15.17. No group of order $56 = 2^3 \cdot 7$ is simple. We have seen that if we can show that there is only one Sylow p-subgroup for some prime p dividing 56, then this must be a normal subgroup and we are done. By the Third Sylow Theorem, there are either one or eight Sylow 7-subgroups. If there is only a single Sylow 7-subgroup, then it must be normal.

On the other hand, suppose that there are eight Sylow 7-subgroups. Then each of these subgroups must be cyclic; hence, the intersection of any two of these subgroups contains only the identity of the group. This leaves $8 \cdot 6 = 48$ distinct elements in the group, each of order 7. Now let us count Sylow 2-subgroups. There are either one or seven Sylow 2-subgroups. Any element of a Sylow 2-subgroup other than the identity must have as its order a power of 2; and therefore cannot be one of the 48 elements of order 7 in the Sylow 7-subgroups. Since a Sylow 2-subgroup has order 8, there is only enough room for a single Sylow 2-subgroup in a group of order 56. If there is only one Sylow 2-subgroup, it must be normal.

For other groups G, it is more difficult to prove that G is not simple. Suppose G has order 48. In this case the technique that we employed in the last example will not work. We need the following lemma to prove that no group of order 48 is simple.

Lemma 15.18. *Let H and K be finite subgroups of a group G. Then*

$$|HK| = \frac{|H| \cdot |K|}{|H \cap K|}.$$

PROOF. Recall that

$$HK = \{hk : h \in H, k \in K\}.$$

Certainly, $|HK| \leq |H| \cdot |K|$ since some element in HK could be written as the product of different elements in H and K. It is quite possible that $h_1 k_1 = h_2 k_2$ for $h_1, h_2 \in H$ and $k_1, k_2 \in K$. If this is the case, let

$$a = (h_1)^{-1} h_2 = k_1 (k_2)^{-1}.$$

Notice that $a \in H \cap K$, since $(h_1)^{-1} h_2$ is in H and $k_2 (k_1)^{-1}$ is in K; consequently,

$$h_2 = h_1 a^{-1}$$
$$k_2 = a k_1.$$

Conversely, let $h = h_1 b^{-1}$ and $k = b k_1$ for $b \in H \cap K$. Then $hk = h_1 k_1$, where $h \in H$ and $k \in K$. Hence, any element $hk \in HK$ can be written in the form $h_i k_i$ for $h_i \in H$ and $k_i \in K$, as many times as there are elements in $H \cap K$; that is, $|H \cap K|$ times. Therefore, $|HK| = (|H| \cdot |K|)/|H \cap K|$. □

Example 15.19. To demonstrate that a group G of order 48 is not simple, we will show that G contains either a normal subgroup of order 8 or a normal subgroup of order 16. By the Third Sylow Theorem, G has either one or three Sylow 2-subgroups of order 16. If there is only one subgroup, then it must be a normal subgroup.

Suppose that the other case is true, and two of the three Sylow 2-subgroups are H and K. We claim that $|H \cap K| = 8$. If $|H \cap K| \leq 4$, then by Lemma 15.18,

$$|HK| = \frac{16 \cdot 16}{4} = 64,$$

which is impossible. Notice that $H \cap K$ has index two in both of H and K, so is normal in both, and thus H and K are each in the normalizer of $H \cap K$. Because H is a subgroup of $N(H \cap K)$ and because $N(H \cap K)$ has strictly more than 16 elements, $|N(H \cap K)|$ must be a multiple of 16 greater than 1, as well as dividing 48. The only possibility is that $|N(H \cap K)| = 48$. Hence, $N(H \cap K) = G$.

The following famous conjecture of Burnside was proved in a long and difficult paper by Feit and Thompson [2].

Theorem 15.20 Odd Order Theorem. *Every finite simple group of nonprime order must be of even order.*

The proof of this theorem laid the groundwork for a program in the 1960s and 1970s that classified all finite simple groups. The success of this program is one of the outstanding achievements of modern mathematics.

15.3   Exercises

1.   What are the orders of all Sylow p-subgroups where G has order 18, 24, 54, 72, and 80?

2.   Find all the Sylow 3-subgroups of S_4 and show that they are all conjugate.

3.   Show that every group of order 45 has a normal subgroup of order 9.

4.   Let H be a Sylow p-subgroup of G. Prove that H is the only Sylow p-subgroup of G contained in $N(H)$.

5.   Prove that no group of order 96 is simple.

6.   Prove that no group of order 160 is simple.

7.   If H is a normal subgroup of a finite group G and $|H| = p^k$ for some prime p, show that H is contained in every Sylow p-subgroup of G.

8.   Let G be a group of order $p^2 q^2$, where p and q are distinct primes such that $q \nmid p^2 - 1$ and $p \nmid q^2 - 1$. Prove that G must be abelian. Find a pair of primes for which this is true.

9.   Show that a group of order 33 has only one Sylow 3-subgroup.

10.   Let H be a subgroup of a group G. Prove or disprove that the normalizer of H is normal in G.

11.   Let G be a finite group divisible by a prime p. Prove that if there is only one Sylow p-subgroup in G, it must be a normal subgroup of G.

12.   Let G be a group of order p^r, p prime. Prove that G contains a normal subgroup of order p^{r-1}.

13.   Suppose that G is a finite group of order $p^n k$, where $k < p$. Show that G must contain a normal subgroup.

14.   Let H be a subgroup of a finite group G. Prove that $gN(H)g^{-1} = N(gHg^{-1})$ for any $g \in G$.

15.   Prove that a group of order 108 must have a normal subgroup.

16.   Classify all the groups of order 175 up to isomorphism.

17.   Show that every group of order 255 is cyclic.

18.   Let G have order $p_1^{e_1} \cdots p_n^{e_n}$ and suppose that G has n Sylow p-subgroups $P_1, \ldots, P_n$ where $|P_i| = p_i^{e_i}$. Prove that G is isomorphic to $P_1 \times \cdots \times P_n$.

19.   Let P be a normal Sylow p-subgroup of G. Prove that every inner automorphism of G fixes P.

20.   What is the smallest possible order of a group G such that G is nonabelian and $|G|$ is odd? Can you find such a group?

21. **The Frattini Lemma.**   If H is a normal subgroup of a finite group G and P is a Sylow p-subgroup of H, for each $g \in G$ show that there is an h in H such that $gPg^{-1} = hPh^{-1}$. Also, show that if N is the normalizer of P, then $G = HN$.

22.   Show that if the order of G is $p^n q$, where p and q are primes and $p > q$, then G contains a normal subgroup.

23.   Prove that the number of distinct conjugates of a subgroup H of a finite group G is $[G : N(H)]$.

24.   Prove that a Sylow 2-subgroup of S_5 is isomorphic to D_4.

25. **Another Proof of the Sylow Theorems.**

 (a) Suppose p is prime and p does not divide m. Show that

$$p \nmid \binom{p^k m}{p^k}.$$

(b) Let $\mathcal{S}$ denote the set of all p^k element subsets of G. Show that p does not divide $|\mathcal{S}|$.

(c) Define an action of G on $\mathcal{S}$ by left multiplication, $aT = \{at : t \in T\}$ for $a \in G$ and $T \in \mathcal{S}$. Prove that this is a group action.

(d) Prove $p \nmid |\mathcal{O}_T|$ for some $T \in \mathcal{S}$.

(e) Let $\{T_1, \ldots, T_u\}$ be an orbit such that $p \nmid u$ and $H = \{g \in G : gT_1 = T_1\}$. Prove that H is a subgroup of G and show that $|G| = u|H|$.

(f) Show that p^k divides $|H|$ and $p^k \leq |H|$.

(g) Show that $|H| = |\mathcal{O}_T| \leq p^k$; conclude that therefore $p^k = |H|$.

26. Let G be a group. Prove that $G' = \langle aba^{-1}b^{-1} : a, b \in G \rangle$ is a normal subgroup of G and G/G' is abelian. Find an example to show that $\{aba^{-1}b^{-1} : a, b \in G\}$ is not necessarily a group.

15.4 A Project

The main objective of finite group theory is to classify all possible finite groups up to isomorphism. This problem is very difficult even if we try to classify the groups of order less than or equal to 60. However, we can break the problem down into several intermediate problems. This is a challenging project that requires a working knowledge of the group theory you have learned up to this point. Even if you do not complete it, it will teach you a great deal about finite groups. You can use Table 15.21 as a guide.

Order	Number	Order	Number	Order	Number	Order	Number
1	?	16	14	31	1	46	2
2	?	17	1	32	51	47	1
3	?	18	?	33	1	48	52
4	?	19	?	34	?	49	?
5	?	20	5	35	1	50	5
6	?	21	?	36	14	51	?
7	?	22	2	37	1	52	?
8	?	23	1	38	?	53	?
9	?	24	?	39	2	54	15
10	?	25	2	40	14	55	2
11	?	26	2	41	1	56	?
12	5	27	5	42	?	57	2
13	?	28	?	43	1	58	?
14	?	29	1	44	4	59	1
15	1	30	4	45	?	60	13

Table 15.21: Numbers of distinct groups G, $|G| \leq 60$

1. Find all simple groups G ($|G| \leq 60$). *Do not use the Odd Order Theorem unless you are prepared to prove it.*

2. Find the number of distinct groups G, where the order of G is n for $n = 1, \ldots, 60$.

3. Find the actual groups (up to isomorphism) for each n.

References and Suggested Readings

[1]   Edwards, H. "A Short History of the Fields Medal," *Mathematical Intelligencer* **1**(1978), 127–29.

[2]   Feit, W. and Thompson, J. G. "Solvability of Groups of Odd Order," *Pacific Journal of Mathematics* **13**(1963), 775–1029.

[3]   Gallian, J. A. "The Search for Finite Simple Groups," *Mathematics Magazine* **49**(1976), 163–79.

[4]   Gorenstein, D. "Classifying the Finite Simple Groups," *Bulletin of the American Mathematical Society* **14**(1986), 1–98.

[5]   Gorenstein, D. *Finite Groups.* AMS Chelsea Publishing, Providence RI, 1968.

[6]   Gorenstein, D., Lyons, R., and Solomon, R. *The Classification of Finite Simple Groups.* American Mathematical Society, Providence RI, 1994.

15.5   Sage

Sylow Subgroups

The Sage permutation group method `.sylow_subgroup(p)` will return a single Sylow p-subgroup. If the prime is not a proper divisor of the group order it returns a subgroup of order p^0, in other words, a trivial subgroup. So be careful about how you construct your primes. Sometimes, you may only want *one* such Sylow subgroup, since any two Sylow p-subgroups are conjugate, and hence isomorphic (Theorem 15.7). This also means we can create other Sylow p-subgroups by conjugating the one we have. The permutation group method `.conjugate(g)` will conjugate the group by `g`.

With repeated conjugations of a single Sylow p-subgroup, we will always create duplicate subgroups. So we need to use a slightly complicated construction to form a list of just the unique subgroups as the list of conjugates. This routine that computes all Sylow p-subgroups can be helpful throughout this section. It could be made much more efficient by conjugating by just one element per coset of the normalizer, but it will be sufficient for our purposes here. Be sure to execute the next cell if you are online, so the function is defined for use later.

```
def all_sylow(G, p):
    '''Form the set of all distinct Sylow p-subgroups of G'''
    scriptP = []
    P = G.sylow_subgroup(p)
    for x in G:
        H = P.conjugate(x)
        if not(H in scriptP):
            scriptP.append(H)
    return scriptP
```

Lets investigate the Sylow subgroups of the dihedral group D_{18}. As a group of order $36 = 2^2 \cdot 3^2$, we know by the First Sylow Theorem that there is a Sylow 2-subgroup of order 4 and a Sylow 3-subgroup of order 9. First for $p = 2$, we obtain one Sylow 2-subgroup, form all the conjugates, and form a list of non-duplicate subgroups. (These commands take a while to execute, so be patient.)

```
G = DihedralGroup(18)
S2 = G.sylow_subgroup(2); S2
```

```
Subgroup of (Dihedral group of order 36 as a permutation group)
generated by
[(2,18)(3,17)(4,16)(5,15)(6,14)(7,13)(8,12)(9,11),
(1,10)(2,11)(3,12)(4,13)(5,14)(6,15)(7,16)(8,17)(9,18)]
```

```
uniqS2 = all_sylow(G, 2)
uniqS2
```

```
[Permutation Group with generators
 [(2,18)(3,17)(4,16)(5,15)(6,14)(7,13)(8,12)(9,11),
 (1,10)(2,11)(3,12)(4,13)(5,14)(6,15)(7,16)(8,17)(9,18)],
 Permutation Group with generators
 [(1,3)(4,18)(5,17)(6,16)(7,15)(8,14)(9,13)(10,12),
 (1,10)(2,11)(3,12)(4,13)(5,14)(6,15)(7,16)(8,17)(9,18)],
 Permutation Group with generators
 [(1,10)(2,11)(3,12)(4,13)(5,14)(6,15)(7,16)(8,17)(9,18),
 (1,17)(2,16)(3,15)(4,14)(5,13)(6,12)(7,11)(8,10)],
 Permutation Group with generators
 [(1,5)(2,4)(6,18)(7,17)(8,16)(9,15)(10,14)(11,13),
 (1,10)(2,11)(3,12)(4,13)(5,14)(6,15)(7,16)(8,17)(9,18)],
 Permutation Group with generators
 [(1,10)(2,11)(3,12)(4,13)(5,14)(6,15)(7,16)(8,17)(9,18),
 (1,15)(2,14)(3,13)(4,12)(5,11)(6,10)(7,9)(16,18)],
 Permutation Group with generators
 [(1,10)(2,11)(3,12)(4,13)(5,14)(6,15)(7,16)(8,17)(9,18),
 (1,13)(2,12)(3,11)(4,10)(5,9)(6,8)(14,18)(15,17)],
 Permutation Group with generators
 [(1,7)(2,6)(3,5)(8,18)(9,17)(10,16)(11,15)(12,14),
 (1,10)(2,11)(3,12)(4,13)(5,14)(6,15)(7,16)(8,17)(9,18)],
 Permutation Group with generators
 [(1,10)(2,11)(3,12)(4,13)(5,14)(6,15)(7,16)(8,17)(9,18),
 (1,11)(2,10)(3,9)(4,8)(5,7)(12,18)(13,17)(14,16)],
 Permutation Group with generators
 [(1,9)(2,8)(3,7)(4,6)(10,18)(11,17)(12,16)(13,15),
 (1,10)(2,11)(3,12)(4,13)(5,14)(6,15)(7,16)(8,17)(9,18)]]
```

```
len(uniqS2)
```

```
9
```

The Third Sylow Theorem tells us that for $p = 2$ we would expect $1, 3$ or 9 Sylow 2-subgroups, so our computational result of 9 subgroups is consistent with what the theory predicts. Can you visualize each of these subgroups as symmetries of an 18-gon? Notice that we also have many subgroups of order 2 inside of these subgroups of order 4.

Now for the case of $p = 3$.

```
G = DihedralGroup(18)
S3 = G.sylow_subgroup(3); S3
```

```
Subgroup of (Dihedral group of order 36 as a permutation group)
generated by
[(1,7,13)(2,8,14)(3,9,15)(4,10,16)(5,11,17)(6,12,18),
(1,15,11,7,3,17,13,9,5)(2,16,12,8,4,18,14,10,6)]
```

```
uniqS3 = all_sylow(G, 3)
uniqS3
```

```
[Permutation Group with generators
[(1,7,13)(2,8,14)(3,9,15)(4,10,16)(5,11,17)(6,12,18),
(1,15,11,7,3,17,13,9,5)(2,16,12,8,4,18,14,10,6)]]
```

```
len(uniqS3)
```

```
1
```

What does the Third Sylow Theorem predict? Just 1 or 4 Sylow 3-subgroups. Having found just one subgroup computationally, we know that all of the conjugates of the lone Sylow 3-subgroup are equal. In other words, the Sylow 3-subgroup is normal in D_{18}. Let us check anyway.

```
S3.is_normal(G)
```

```
True
```

At least one of the subgroups of order 3 contained in this Sylow 3-subgroup should be obvious by looking at the orders of the generators, and then you may even notice that the generators given could be reduced, and one is a power of the other.

```
S3.is_cyclic()
```

```
True
```

Remember that there are many other subgroups, of other orders. For example, can you construct a subgroup of order $6 = 2 \cdot 3$ in D_{18}?

Normalizers

A new command that is relevant to this section is the construction of a normalizer. The Sage command `G.normalizer(H)` will return the subgroup of `G` containing elements that normalize the subgroup `H`. We illustrate its use with the Sylow subgroups from above.

```
G = DihedralGroup(18)
S2 = G.sylow_subgroup(2)
S3 = G.sylow_subgroup(3)
N2 = G.normalizer(S2); N2
```

```
Subgroup of (Dihedral group of order 36 as a permutation group)
generated by
[(2,18)(3,17)(4,16)(5,15)(6,14)(7,13)(8,12)(9,11),
(1,10)(2,11)(3,12)(4,13)(5,14)(6,15)(7,16)(8,17)(9,18)]
```

```
N2 == S2
```

```
True
```

```
N3 = G.normalizer(S3); N3
```

```
Subgroup of (Dihedral group of order 36 as a permutation group)
generated by
[(2,18)(3,17)(4,16)(5,15)(6,14)(7,13)(8,12)(9,11),
(1,2)(3,18)(4,17)(5,16)(6,15)(7,14)(8,13)(9,12)(10,11),
(1,7,13)(2,8,14)(3,9,15)(4,10,16)(5,11,17)(6,12,18),
(1,15,11,7,3,17,13,9,5)(2,16,12,8,4,18,14,10,6)]
```

```
N3 == G
```

```
True
```

The normalizer of a subgroup always contains the whole subgroup, so the normalizer of S2 is as small as possible. We already knew S3 is normal in G, so it is no surprise that its normalizer is as big as possible — every element of G normalizes S3. Let us compute a normalizer in D_{18} that is more "interesting."

```
G = DihedralGroup(18)
a = G("(1,7,13)(2,8,14)(3,9,15)(4,10,16)(5,11,17)(6,12,18)")
b = G("(1,5)(2,4)(6,18)(7,17)(8,16)(9,15)(10,14)(11,13)")
H = G.subgroup([a, b])
H.order()
```

```
6
```

```
N = G.normalizer(H)
N
```

```
Subgroup of (Dihedral group of order 36 as a permutation group)
generated by
[(1,2)(3,18)(4,17)(5,16)(6,15)(7,14)(8,13)(9,12)(10,11),
(1,5)(2,4)(6,18)(7,17)(8,16)(9,15)(10,14)(11,13),
(1,13,7)(2,14,8)(3,15,9)(4,16,10)(5,17,11)(6,18,12)]
```

```
N.order()
```

```
12
```

So for this subgroup of order 6, the normalizer is strictly bigger than the subgroup, but still strictly smaller than the whole group (and hence not normal in the dihedral group). Trivially, a subgroup is normal in its normalizer:

```
H.is_normal(G)
```

```
False
```

```
H.is_normal(N)
```

```
True
```

Finite Simple Groups

We saw earlier Sage's permutation group method `.is_simple()`. Example 15.16 tells us that a group of order 64 is never simple. The dicyclic group `DiCyclicGroup(16)` is a non-abelian group of 64, so we can test this method on this group. It turns out this group has many normal subgroups — the list will always contain the trivial subgroup and the group itself, so any number exceeding 2 indicates a non-trivial normal subgroup.

```
DC=DiCyclicGroup(16)
DC.order()
```

```
64
```

```
DC.is_simple()
```

```
False
```

```
ns = DC.normal_subgroups()
len(ns)
```

9

Here is a rather interesting group, one of the 26 sporadic simple groups, known as the Higman-Sims group, HS. The generators used below come from the representation on 100 points in GAP format, available off of web.mat.bham.ac.uk/atlas/v2.0/spor/HS/. Two generators of just order 2 and order 5 (as you can esily see), generating exactly $44\,352\,000$ elements, but no normal subgroups. Amazing.

```
G = SymmetricGroup(100)
a = G([(1,60),  (2,72),  (3,81),  (4,43),  (5,11),  (6,87),
       (7,34),  (9,63),  (12,46), (13,28), (14,71), (15,42),
       (16,97), (18,57), (19,52), (21,32), (23,47), (24,54),
       (25,83), (26,78), (29,89), (30,39), (33,61), (35,56),
       (37,67), (44,76), (45,88), (48,59), (49,86), (50,74),
       (51,66), (53,99), (55,75), (62,73), (65,79), (68,82),
       (77,92), (84,90), (85,98), (94,100)])
b = G([(1,86,13,10,47),  (2,53,30,8,38),
       (3,40,48,25,17),  (4,29,92,88,43),  (5,98,66,54, 65),
       (6,27,51,73,24),  (7,83,16,20,28),  (9,23,89,95,61),
       (11,42,46,91,32), (12,14, 81,55,68), (15,90,31,56,37),
       (18,69,45,84,76), (19,59,79,35,93),  (21,22,64,39,100),
       (26,58,96,85,77), (33,52,94,75,44),  (34,62,87,78,50),
       (36,82,60,74,72), (41,80,70,49,67),  (57,63,71,99,97)])
a.order(), b.order()
```

(2, 5)

```
HS = G.subgroup([a, b])
HS.order()
```

44352000

```
HS.is_simple()
```

True

We saw this group earlier in the exercises for Chapter 14 on group actions, where it was the single non-trivial normal subgroup of the automorphism group of the Higman-Sims graph, hence its name.

GAP Console and Interface

This concludes our exclusive study of group theory, though we will be using groups some in the subsequent sections. As we have remarked, much of Sage's computation with groups is performed by the open source program, "Groups, Algorithms, and Programming," which is better know as simply GAP. If after this course you outgrow Sage's support for groups, then learning GAP would be your next step as a group theorist. Every copy of Sage includes a copy of GAP and is easy to see which version of GAP is included:

```
gap.version()
```

'4.8.6'

You can interact with GAP in Sage in several ways. The most direct is by creating a permutation group via Sage's `gap()` command.

```
G = gap('Group(␣(1,2,3,4,5,6),␣(1,3,5)␣)')
G
```

```
Group( [ (1,2,3,4,5,6), (1,3,5) ] )
```

Now we can use most any GAP command with `G`, via the convention that most GAP commands expect a group as the first argument, and we instead provide the group by using the object-oriented`G`. syntax. If you consult the GAP documentation you will see that `Center` is a GAP command that expects a group as its lone argument, and `Centralizer` is a GAP command that expects two arguments — a group and then a group element.

```
G.Center()
```

```
Group( [ ( 1, 3, 5)( 2, 4, 6) ] )
```

```
G.Centralizer('(1,␣3,␣5)')
```

```
Group( [ (1,3,5), (2,4,6), (1,3,5)(2,4,6) ] )
```

If you use the Sage Notebook interface you can set the first line of a compute cell to `%gap` and the entire cell will be interpreted as if you were interacting directly with GAP. This means you would now use GAP's syntax, which you can see above is slightly different than Sage's universal syntax. You can also use the drop-down box at the top of a worksheet, and select `gap` as the system (rather than `sage`) and your whole worksheet will be interpreted as GAP commands. Here is one simple example, which you should be able to evaluate in your current worksheet. This particular example will not run properly in a Sage Cell in a web page version of this section.

```
%gap
G := Group( (1,2,3,4,5,6), (1,3,5) );
Centralizer(G, (1,3,5));
```

Notice that

- We do not need to wrap the individual permutations in as many quotation marks as we do in Sage.

- Assignment is `:=` not `=`. If you forget the colon, you will get an error message such as `Variable: 'G' must have a value`

- A line *must* end with a semi-colon. If you forget, several lines will be merged together.

You can get help about GAP commands with a command such as the following, though you will soon see that GAP assumes you know a lot more algebra than Sage assumes you know.

```
print(gap.help('SymmetricGroup', pager=False))
```

In the command-line version of Sage, you can also use the GAP "console." Again, you need to use GAP syntax, and you do not have many of the conveniences of the Sage notebook. It is also good to know in advance that `quit;` is how you can leave the GAP console and get back to Sage. If you run Sage at the command-line, use the command `gap_console()` to start GAP running.

It is a comfort to know that with Sage you get a complete copy of GAP, installed and all ready to run. However, this is not a tutorial on GAP, so consult the documentation available at the main GAP website: www.gap-system.org to learn how to get the most out of GAP.

15.6   Sage Exercises

1.   This exercise verifies Theorem 15.13. The commutator subgroup is computed with the permutation group method `.commutator()`. For the dihedral group of order 40, D_{20} (`DihedralGroup(20)` in Sage), compute the commutator subgroup and form the quotient with the dihedral group. Then verify that this quotient is abelian. Can you identify the quotient group exactly (in other words, up to isomorphism)?

2.   For each possible prime, find all of the distinct Sylow p-subgroups of the alternating group A_5. Confirm that your results are consistent with the Third Sylow Theorem for each prime. We know that A_5 is a simple group. Explain how this would explain or predict some aspects of your answers.

   Count the number of distinct elements contained in the union of all the Sylow subgroups you just found. What is interesting about this count?

3.   For the dihedral group D_{36} (symmetries of a 36-gon) and each possible prime, determine the possibilities for the number of distinct Sylow p-subgroups as predicted by the Third Sylow Theorem (15.8). Now compute the actual number of distinct Sylow p-subgroups for each prime and comment on the result.

   It can be proved that *any group* with order 72 is not a simple group, using techniques such as those used in the later examples in this chapter. Discuss this result in the context of your computations with Sage.

4.   This exercise verifies Lemma 15.6. Let G be the dihedral group of order 36, D_{18}. Let H be the one Sylow 3-subgroup. Let K be the subgroup of order 6 generated by the two permutations `a` and `b` given below. First, form a list of the distinct conjugates of K by the elements of H, and determine the number of subgroups in this list. Compare this with the index given in the statement of the lemma, employing a single (long) statement making use of the `.order()`, `.normalizer()` and `.intersection()` methods with G, H and K, *only*.

```
G = DihedralGroup(18)
a = G("(1,7,13)(2,8,14)(3,9,15)(4,10,16)(5,11,17)(6,12,18)")
b = G("(1,5)(2,4)(6,18)(7,17)(8,16)(9,15)(10,14)(11,13)")
```

5.   Example 15.19 shows that every group of order 48 has a normal subgroup. The dicyclic groups are an infinite family of non-abelian groups with order $4n$, which includes the quaternions (the case of $n = 2$). So the permutation group `DiCyclicGroup(12)` has order 48. Use Sage to follow the logic of the proof in Example 15.19 and construct a normal subgroup in this group. (In other words, do not just ask for a list of the normal subgroups from Sage, but instead trace through the implications in the example to arrive at a normal subgroup, and then check your answer.)

6.   The proofs of the Second and Third Sylow Theorems (15.7, 15.8) employ a group action on sets of Sylow p-subgroups, For the Second Theorem, the list is proposed as incomplete and is proved to be *all* of the Sylow p-subgroups. In this exercise we will see how these actions behave, and how they are different when we use different groups acting on the same set.

   Construct the six Sylow 5-subgroups of the alternating group A_5. This will be the set of objects for both of our actions. Conjugating one of these Sylow 5-subgroups by an element of A_5 will produce another Sylow 5-subgroup, and so can be used to create a group action. For such an action, from each group element form a Sage permutation of the subgroups by numbering the six subgroups and using these integers as markers for the subgroups. You will find the Python list method `.index()` very helpful. Now use all of these permutations to generate a permutation group (a subgroup of S_6). Finally, use permutation group methods

for orbits and stabilisers, etc. to explore the actions.

For the first action, use all of A_5 as the group. Show that the resulting action is transitive. In other words, there is exactly one single orbit.

For the second action, use just one of the Sylow 5-subgroups as the group. Write the class equation for this action in a format that suggests the "congruent to 1 mod p" part of the conclusion of the Third Theorem.

Rings

Up to this point we have studied sets with a single binary operation satisfying certain axioms, but we are often more interested in working with sets that have two binary operations. For example, one of the most natural algebraic structures to study is the integers with the operations of addition and multiplication. These operations are related to one another by the distributive property. If we consider a set with two such related binary operations satisfying certain axioms, we have an algebraic structure called a ring. In a ring we add and multiply elements such as real numbers, complex numbers, matrices, and functions.

16.1 Rings

A nonempty set R is a ***ring*** if it has two closed binary operations, addition and multiplication, satisfying the following conditions.

1. $a + b = b + a$ for $a, b \in R$.

2. $(a + b) + c = a + (b + c)$ for $a, b, c \in R$.

3. There is an element 0 in R such that $a + 0 = a$ for all $a \in R$.

4. For every element $a \in R$, there exists an element $-a$ in R such that $a + (-a) = 0$.

5. $(ab)c = a(bc)$ for $a, b, c \in R$.

6. For $a, b, c \in R$,

$$a(b + c) = ab + ac$$
$$(a + b)c = ac + bc.$$

This last condition, the distributive axiom, relates the binary operations of addition and multiplication. Notice that the first four axioms simply require that a ring be an abelian group under addition, so we could also have defined a ring to be an abelian group $(R, +)$ together with a second binary operation satisfying the fifth and sixth conditions given above.

If there is an element $1 \in R$ such that $1 \neq 0$ and $1a = a1 = a$ for each element $a \in R$, we say that R is a ring with ***unity*** or ***identity***. A ring R for which $ab = ba$ for all a, b in R is called a ***commutative ring***. A commutative ring R with identity is called an ***integral domain*** if, for every $a, b \in R$ such that $ab = 0$, either $a = 0$ or $b = 0$. A ***division ring*** is a ring R, with an identity, in which every nonzero element in R is a ***unit***; that is, for each $a \in R$ with $a \neq 0$, there exists a unique element a^{-1} such that $a^{-1}a = aa^{-1} = 1$. A commutative division ring is called a ***field***. The relationship among rings, integral domains, division rings, and fields is shown in Figure 16.1.
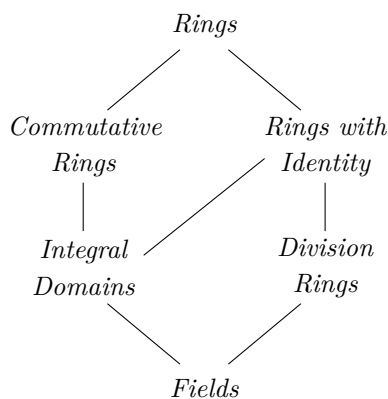
Figure 16.1: Types of rings

Example 16.2. As we have mentioned previously, the integers form a ring. In fact, $\mathbb{Z}$ is an integral domain. Certainly if $ab = 0$ for two integers a and b, either $a = 0$ or $b = 0$. However, $\mathbb{Z}$ is not a field. There is no integer that is the multiplicative inverse of 2, since $1/2$ is not an integer. The only integers with multiplicative inverses are 1 and -1.

Example 16.3. Under the ordinary operations of addition and multiplication, all of the familiar number systems are rings: the rationals, $\mathbb{Q}$; the real numbers, $\mathbb{R}$; and the complex numbers, $\mathbb{C}$. Each of these rings is a field.

Example 16.4. We can define the product of two elements a and b in $\mathbb{Z}_n$ by $ab \pmod{n}$. For instance, in $\mathbb{Z}_{12}$, $5 \cdot 7 \equiv 11 \pmod{12}$. This product makes the abelian group $\mathbb{Z}_n$ into a ring. Certainly $\mathbb{Z}_n$ is a commutative ring; however, it may fail to be an integral domain. If we consider $3 \cdot 4 \equiv 0 \pmod{12}$ in $\mathbb{Z}_{12}$, it is easy to see that a product of two nonzero elements in the ring can be equal to zero.

A nonzero element a in a ring R is called a ***zero divisor*** if there is a nonzero element b in R such that $ab = 0$. In the previous example, 3 and 4 are zero divisors in $\mathbb{Z}_{12}$.

Example 16.5. In calculus the continuous real-valued functions on an interval $[a, b]$ form a commutative ring. We add or multiply two functions by adding or multiplying the values of the functions. If $f(x) = x^2$ and $g(x) = \cos x$, then $(f + g)(x) = f(x) + g(x) = x^2 + \cos x$ and $(fg)(x) = f(x)g(x) = x^2 \cos x$.

Example 16.6. The 2×2 matrices with entries in $\mathbb{R}$ form a ring under the usual operations of matrix addition and multiplication. This ring is noncommutative, since it is usually the case that $AB \neq BA$. Also, notice that we can have $AB = 0$ when neither A nor B is zero.

Example 16.7. For an example of a noncommutative division ring, let

$$1 = \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}, \quad \mathbf{i} = \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}, \quad \mathbf{j} = \begin{pmatrix} 0 & i \\ i & 0 \end{pmatrix}, \quad \mathbf{k} = \begin{pmatrix} i & 0 \\ 0 & -i \end{pmatrix},$$

where $i^2 = -1$. These elements satisfy the following relations:

$$\mathbf{i}^2 = \mathbf{j}^2 = \mathbf{k}^2 = -1$$
$$\mathbf{ij} = \mathbf{k}$$
$$\mathbf{jk} = \mathbf{i}$$
$$\mathbf{ki} = \mathbf{j}$$
$$\mathbf{ji} = -\mathbf{k}$$

$$\mathbf{kj} = -\mathbf{i}$$
$$\mathbf{ik} = -\mathbf{j}.$$

Let $\mathbb{H}$ consist of elements of the form $a + b\mathbf{i} + c\mathbf{j} + d\mathbf{k}$, where a, b, c, d are real numbers. Equivalently, $\mathbb{H}$ can be considered to be the set of all 2×2 matrices of the form

$$\begin{pmatrix} \alpha & \beta \\ -\overline{\beta} & \overline{\alpha} \end{pmatrix},$$

where $\alpha = a + di$ and $\beta = b + ci$ are complex numbers. We can define addition and multiplication on $\mathbb{H}$ either by the usual matrix operations or in terms of the generators 1, $\mathbf{i}$, $\mathbf{j}$, and $\mathbf{k}$:

$$(a_1 + b_1\mathbf{i} + c_1\mathbf{j} + d_1\mathbf{k}) + (a_2 + b_2\mathbf{i} + c_2\mathbf{j} + d_2\mathbf{k})$$
$$= (a_1 + a_2) + (b_1 + b_2)\mathbf{i} + (c_1 + c_2)\mathbf{j} + (d_1 + d_2)\mathbf{k}$$

and

$$(a_1 + b_1\mathbf{i} + c_1\mathbf{j} + d_1\mathbf{k})(a_2 + b_2\mathbf{i} + c_2\mathbf{j} + d_2\mathbf{k}) = \alpha + \beta\mathbf{i} + \gamma\mathbf{j} + \delta\mathbf{k},$$

where

$$\alpha = a_1 a_2 - b_1 b_2 - c_1 c_2 - d_1 d_2$$
$$\beta = a_1 b_2 + a_2 b_1 + c_1 d_2 - d_1 c_2$$
$$\gamma = a_1 c_2 - b_1 d_2 + c_1 a_2 + d_1 b_2$$
$$\delta = a_1 d_2 + b_1 c_2 - c_1 b_2 + d_1 a_2.$$

Though multiplication looks complicated, it is actually a straightforward computation if we remember that we just add and multiply elements in $\mathbb{H}$ like polynomials and keep in mind the relationships between the generators $\mathbf{i}$, $\mathbf{j}$, and $\mathbf{k}$. The ring $\mathbb{H}$ is called the ring of **quaternions**.

To show that the quaternions are a division ring, we must be able to find an inverse for each nonzero element. Notice that

$$(a + b\mathbf{i} + c\mathbf{j} + d\mathbf{k})(a - b\mathbf{i} - c\mathbf{j} - d\mathbf{k}) = a^2 + b^2 + c^2 + d^2.$$

This element can be zero only if a, b, c, and d are all zero. So if $a + b\mathbf{i} + c\mathbf{j} + d\mathbf{k} \neq 0$,

$$(a + b\mathbf{i} + c\mathbf{j} + d\mathbf{k})\left(\frac{a - b\mathbf{i} - c\mathbf{j} - d\mathbf{k}}{a^2 + b^2 + c^2 + d^2}\right) = 1.$$

Proposition 16.8. *Let R be a ring with $a, b \in R$. Then*

1. *$a0 = 0a = 0$;*

2. *$a(-b) = (-a)b = -ab$;*

3. *$(-a)(-b) = ab$.*

PROOF.   To prove (1), observe that

$$a0 = a(0 + 0) = a0 + a0;$$

hence, $a0 = 0$. Similarly, $0a = 0$. For (2), we have $ab + a(-b) = a(b - b) = a0 = 0$; consequently, $-ab = a(-b)$. Similarly, $-ab = (-a)b$. Part (3) follows directly from (2) since $(-a)(-b) = -(a(-b)) = -(-ab) = ab$.   □

Just as we have subgroups of groups, we have an analogous class of substructures for rings. A **subring** S of a ring R is a subset S of R such that S is also a ring under the inherited operations from R.

Example 16.9. The ring $n\mathbb{Z}$ is a subring of $\mathbb{Z}$. Notice that even though the original ring may have an identity, we do not require that its subring have an identity. We have the following chain of subrings:

$$\mathbb{Z} \subset \mathbb{Q} \subset \mathbb{R} \subset \mathbb{C}.$$

The following proposition gives us some easy criteria for determining whether or not a subset of a ring is indeed a subring. (We will leave the proof of this proposition as an exercise.)

Proposition 16.10. *Let R be a ring and S a subset of R. Then S is a subring of R if and only if the following conditions are satisfied.*

1. $S \neq \emptyset$.

2. $rs \in S$ for all $r, s \in S$.

3. $r - s \in S$ for all $r, s \in S$.

Example 16.11. Let $R = \mathbb{M}_2(\mathbb{R})$ be the ring of 2×2 matrices with entries in $\mathbb{R}$. If T is the set of upper triangular matrices in R; i.e.,

$$T = \left\{ \begin{pmatrix} a & b \\ 0 & c \end{pmatrix} : a, b, c \in \mathbb{R} \right\},$$

then T is a subring of R. If

$$A = \begin{pmatrix} a & b \\ 0 & c \end{pmatrix} \quad \text{and} \quad B = \begin{pmatrix} a' & b' \\ 0 & c' \end{pmatrix}$$

are in T, then clearly $A - B$ is also in T. Also,

$$AB = \begin{pmatrix} aa' & ab' + bc' \\ 0 & cc' \end{pmatrix}$$

is in T.

16.2 Integral Domains and Fields

Let us briefly recall some definitions. If R is a ring and r is a nonzero element in R, then r is said to be a **zero divisor** if there is some nonzero element $s \in R$ such that $rs = 0$. A commutative ring with identity is said to be an **integral domain** if it has no zero divisors. If an element a in a ring R with identity has a multiplicative inverse, we say that a is a **unit**. If every nonzero element in a ring R is a unit, then R is called a **division ring**. A commutative division ring is called a **field**.

Example 16.12. If $i^2 = -1$, then the set $\mathbb{Z}[i] = \{m + ni : m, n \in \mathbb{Z}\}$ forms a ring known as the **Gaussian integers**. It is easily seen that the Gaussian integers are a subring of the complex numbers since they are closed under addition and multiplication. Let $\alpha = a + bi$ be a unit in $\mathbb{Z}[i]$. Then $\overline{\alpha} = a - bi$ is also a unit since if $\alpha\beta = 1$, then $\overline{\alpha}\overline{\beta} = 1$. If $\beta = c + di$, then

$$1 = \alpha\beta\overline{\alpha}\overline{\beta} = (a^2 + b^2)(c^2 + d^2).$$

Therefore, $a^2 + b^2$ must either be 1 or -1; or, equivalently, $a + bi = \pm 1$ or $a + bi = \pm i$. Therefore, units of this ring are ± 1 and $\pm i$; hence, the Gaussian integers are not a field. We will leave it as an exercise to prove that the Gaussian integers are an integral domain.

Example 16.13. The set of matrices

$$F = \left\{ \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}, \begin{pmatrix} 1 & 1 \\ 1 & 0 \end{pmatrix}, \begin{pmatrix} 0 & 1 \\ 1 & 1 \end{pmatrix}, \begin{pmatrix} 0 & 0 \\ 0 & 0 \end{pmatrix} \right\}$$

with entries in $\mathbb{Z}_2$ forms a field.

Example 16.14. The set $\mathbb{Q}(\sqrt{2}) = \{a + b\sqrt{2} : a, b \in \mathbb{Q}\}$ is a field. The inverse of an element $a + b\sqrt{2}$ in $\mathbb{Q}(\sqrt{2})$ is

$$\frac{a}{a^2 - 2b^2} + \frac{-b}{a^2 - 2b^2}\sqrt{2}.$$

We have the following alternative characterization of integral domains.

Proposition 16.15 Cancellation Law. *Let D be a commutative ring with identity. Then D is an integral domain if and only if for all nonzero elements $a \in D$ with $ab = ac$, we have $b = c$.*

PROOF.   Let D be an integral domain. Then D has no zero divisors. Let $ab = ac$ with $a \neq 0$. Then $a(b - c) = 0$. Hence, $b - c = 0$ and $b = c$.

Conversely, let us suppose that cancellation is possible in D. That is, suppose that $ab = ac$ implies $b = c$. Let $ab = 0$. If $a \neq 0$, then $ab = a0$ or $b = 0$. Therefore, a cannot be a zero divisor. $\square$

The following surprising theorem is due to Wedderburn.

Theorem 16.16. *Every finite integral domain is a field.*

PROOF.   Let D be a finite integral domain and D^* be the set of nonzero elements of D. We must show that every element in D^* has an inverse. For each $a \in D^*$ we can define a map $\lambda_a : D^* \to D^*$ by $\lambda_a(d) = ad$. This map makes sense, because if $a \neq 0$ and $d \neq 0$, then $ad \neq 0$. The map λ_a is one-to-one, since for $d_1, d_2 \in D^*$,

$$ad_1 = \lambda_a(d_1) = \lambda_a(d_2) = ad_2$$

implies $d_1 = d_2$ by left cancellation. Since D^* is a finite set, the map λ_a must also be onto; hence, for some $d \in D^*$, $\lambda_a(d) = ad = 1$. Therefore, a has a left inverse. Since D is commutative, d must also be a right inverse for a. Consequently, D is a field. $\square$

For any nonnegative integer n and any element r in a ring R we write $r + \cdots + r$ (n times) as nr. We define the ***characteristic*** of a ring R to be the least positive integer n such that $nr = 0$ for all $r \in R$. If no such integer exists, then the characteristic of R is defined to be 0. We will denote the characteristic of R by char R.

Example 16.17. For every prime p, $\mathbb{Z}_p$ is a field of characteristic p. By Proposition 3.4, every nonzero element in $\mathbb{Z}_p$ has an inverse; hence, $\mathbb{Z}_p$ is a field. If a is any nonzero element in the field, then $pa = 0$, since the order of any nonzero element in the abelian group $\mathbb{Z}_p$ is p.

Lemma 16.18. *Let R be a ring with identity. If 1 has order n, then the characteristic of R is n.*

PROOF.   If 1 has order n, then n is the least positive integer such that $n1 = 0$. Thus, for all $r \in R$,

$$nr = n(1r) = (n1)r = 0r = 0.$$

On the other hand, if no positive n exists such that $n1 = 0$, then the characteristic of R is zero. $\square$

Theorem 16.19. *The characteristic of an integral domain is either prime or zero.*

PROOF.   Let D be an integral domain and suppose that the characteristic of D is n with $n \neq 0$. If n is not prime, then $n = ab$, where $1 < a < n$ and $1 < b < n$. By Lemma 16.18, we need only consider the case $n1 = 0$. Since $0 = n1 = (ab)1 = (a1)(b1)$ and there are no zero divisors in D, either $a1 = 0$ or $b1 = 0$. Hence, the characteristic of D must be less than n, which is a contradiction. Therefore, n must be prime. $\square$

16.3   Ring Homomorphisms and Ideals

In the study of groups, a homomorphism is a map that preserves the operation of the group. Similarly, a homomorphism between rings preserves the operations of addition and multiplication in the ring. More specifically, if R and S are rings, then a **ring homomorphism** is a map $\phi : R \to S$ satisfying

$$\phi(a + b) = \phi(a) + \phi(b)$$
$$\phi(ab) = \phi(a)\phi(b)$$

for all $a, b \in R$. If $\phi : R \to S$ is a one-to-one and onto homomorphism, then ϕ is called an **isomorphism** of rings.

The set of elements that a ring homomorphism maps to 0 plays a fundamental role in the theory of rings. For any ring homomorphism $\phi : R \to S$, we define the **kernel** of a ring homomorphism to be the set

$$\ker \phi = \{r \in R : \phi(r) = 0\}.$$

Example 16.20. For any integer n we can define a ring homomorphism $\phi : \mathbb{Z} \to \mathbb{Z}_n$ by $a \mapsto a \pmod{n}$. This is indeed a ring homomorphism, since

$$\phi(a + b) = (a + b) \pmod{n}$$
$$= a \pmod{n} + b \pmod{n}$$
$$= \phi(a) + \phi(b)$$

and

$$\phi(ab) = ab \pmod{n}$$
$$= a \pmod{n} \cdot b \pmod{n}$$
$$= \phi(a)\phi(b).$$

The kernel of the homomorphism ϕ is $n\mathbb{Z}$.

Example 16.21. Let $C[a, b]$ be the ring of continuous real-valued functions on an interval $[a, b]$ as in Example 16.5. For a fixed $\alpha \in [a, b]$, we can define a ring homomorphism $\phi_\alpha : C[a, b] \to \mathbb{R}$ by $\phi_\alpha(f) = f(\alpha)$. This is a ring homomorphism since

$$\phi_\alpha(f + g) = (f + g)(\alpha) = f(\alpha) + g(\alpha) = \phi_\alpha(f) + \phi_\alpha(g)$$
$$\phi_\alpha(fg) = (fg)(\alpha) = f(\alpha)g(\alpha) = \phi_\alpha(f)\phi_\alpha(g).$$

Ring homomorphisms of the type ϕ_α are called **evaluation homomorphisms**.

In the next proposition we will examine some fundamental properties of ring homomorphisms. The proof of the proposition is left as an exercise.

Proposition 16.22. *Let $\phi : R \to S$ be a ring homomorphism.*

1. *If R is a commutative ring, then $\phi(R)$ is a commutative ring.*

2. *$\phi(0) = 0$.*

3. *Let 1_R and 1_S be the identities for R and S, respectively. If ϕ is onto, then $\phi(1_R) = 1_S$.*

4. *If R is a field and $\phi(R) \neq \{0\}$, then $\phi(R)$ is a field.*

In group theory we found that normal subgroups play a special role. These subgroups have nice characteristics that make them more interesting to study than arbitrary subgroups. In ring theory the objects corresponding to normal subgroups are a special class of subrings called ideals. An ***ideal*** in a ring R is a subring I of R such that if a is in I and r is in R, then both ar and ra are in I; that is, $rI \subset I$ and $Ir \subset I$ for all $r \in R$.

Example 16.23. Every ring R has at least two ideals, $\{0\}$ and R. These ideals are called the ***trivial ideals***.

Let R be a ring with identity and suppose that I is an ideal in R such that 1 is in I. Since for any $r \in R$, $r1 = r \in I$ by the definition of an ideal, $I = R$.

Example 16.24. If a is any element in a commutative ring R with identity, then the set

$$\langle a \rangle = \{ar : r \in R\}$$

is an ideal in R. Certainly, $\langle a \rangle$ is nonempty since both $0 = a0$ and $a = a1$ are in $\langle a \rangle$. The sum of two elements in $\langle a \rangle$ is again in $\langle a \rangle$ since $ar + ar' = a(r + r')$. The inverse of ar is $-ar = a(-r) \in \langle a \rangle$. Finally, if we multiply an element $ar \in \langle a \rangle$ by an arbitrary element $s \in R$, we have $s(ar) = a(sr)$. Therefore, $\langle a \rangle$ satisfies the definition of an ideal.

If R is a commutative ring with identity, then an ideal of the form $\langle a \rangle = \{ar : r \in R\}$ is called a ***principal ideal***.

Theorem 16.25. *Every ideal in the ring of integers $\mathbb{Z}$ is a principal ideal.*

PROOF.   The zero ideal $\{0\}$ is a principal ideal since $\langle 0 \rangle = \{0\}$. If I is any nonzero ideal in $\mathbb{Z}$, then I must contain some positive integer m. There exists a least positive integer n in I by the Principle of Well-Ordering. Now let a be any element in I. Using the division algorithm, we know that there exist integers q and r such that

$$a = nq + r$$

where $0 \leq r < n$. This equation tells us that $r = a - nq \in I$, but r must be 0 since n is the least positive element in I. Therefore, $a = nq$ and $I = \langle n \rangle$. $\square$

Example 16.26. The set $n\mathbb{Z}$ is ideal in the ring of integers. If na is in $n\mathbb{Z}$ and b is in $\mathbb{Z}$, then nab is in $n\mathbb{Z}$ as required. In fact, by Theorem 16.25, these are the only ideals of $\mathbb{Z}$.

Proposition 16.27. *The kernel of any ring homomorphism $\phi : R \to S$ is an ideal in R.*

PROOF.   We know from group theory that $\ker \phi$ is an additive subgroup of R. Suppose that $r \in R$ and $a \in \ker \phi$. Then we must show that ar and ra are in $\ker \phi$. However,

$$\phi(ar) = \phi(a)\phi(r) = 0\phi(r) = 0$$

and

$$\phi(ra) = \phi(r)\phi(a) = \phi(r)0 = 0.$$ $\square$

Remark 16.28. In our definition of an ideal we have required that $rI \subset I$ and $Ir \subset I$ for all $r \in R$. Such ideals are sometimes referred to as **two-sided ideals**. We can also consider **one-sided ideals**; that is, we may require only that either $rI \subset I$ or $Ir \subset I$ for $r \in R$ hold but not both. Such ideals are called **left ideals** and **right ideals**, respectively. Of course, in a commutative ring any ideal must be two-sided. In this text we will concentrate on two-sided ideals.

Theorem 16.29. *Let I be an ideal of R. The factor group R/I is a ring with multiplication defined by*
$$(r + I)(s + I) = rs + I.$$

PROOF. We already know that R/I is an abelian group under addition. Let $r + I$ and $s + I$ be in R/I. We must show that the product $(r + I)(s + I) = rs + I$ is independent of the choice of coset; that is, if $r' \in r + I$ and $s' \in s + I$, then $r's'$ must be in $rs + I$. Since $r' \in r + I$, there exists an element a in I such that $r' = r + a$. Similarly, there exists a $b \in I$ such that $s' = s + b$. Notice that
$$r's' = (r + a)(s + b) = rs + as + rb + ab$$

and $as + rb + ab \in I$ since I is an ideal; consequently, $r's' \in rs + I$. We will leave as an exercise the verification of the associative law for multiplication and the distributive laws. $\square$

The ring R/I in Theorem 16.29 is called the **factor** or **quotient ring**. Just as with group homomorphisms and normal subgroups, there is a relationship between ring homomorphisms and ideals.

Theorem 16.30. *Let I be an ideal of R. The map $\phi : R \to R/I$ defined by $\phi(r) = r + I$ is a ring homomorphism of R onto R/I with kernel I.*

PROOF. Certainly $\phi : R \to R/I$ is a surjective abelian group homomorphism. It remains to show that ϕ works correctly under ring multiplication. Let r and s be in R. Then
$$\phi(r)\phi(s) = (r + I)(s + I) = rs + I = \phi(rs),$$

which completes the proof of the theorem. $\square$

The map $\phi : R \to R/I$ is often called the **natural** or **canonical homomorphism**. In ring theory we have isomorphism theorems relating ideals and ring homomorphisms similar to the isomorphism theorems for groups that relate normal subgroups and homomorphisms in Chapter 11. We will prove only the First Isomorphism Theorem for rings in this chapter and leave the proofs of the other two theorems as exercises. All of the proofs are similar to the proofs of the isomorphism theorems for groups.

Theorem 16.31 First Isomorphism Theorem. *Let $\psi : R \to S$ be a ring homomorphism. Then $\ker \psi$ is an ideal of R. If $\phi : R \to R/\ker \psi$ is the canonical homomorphism, then there exists a unique isomorphism $\eta : R/\ker \psi \to \psi(R)$ such that $\psi = \eta\phi$.*

PROOF. Let $K = \ker \psi$. By the First Isomorphism Theorem for groups, there exists a well-defined group homomorphism $\eta : R/K \to \psi(R)$ defined by $\eta(r + K) = \psi(r)$ for the additive abelian groups R and R/K. To show that this is a ring homomorphism, we need only show that $\eta((r + K)(s + K)) = \eta(r + K)\eta(s + K)$; but
$$\begin{aligned}
\eta((r + K)(s + K)) &= \eta(rs + K) \\
&= \psi(rs) \\
&= \psi(r)\psi(s) \\
&= \eta(r + K)\eta(s + K).
\end{aligned}$$
$\square$

Theorem 16.32 Second Isomorphism Theorem. *Let I be a subring of a ring R and J an ideal of R. Then $I \cap J$ is an ideal of I and*

$$I/I \cap J \cong (I + J)/J.$$

Theorem 16.33 Third Isomorphism Theorem. *Let R be a ring and I and J be ideals of R where $J \subset I$. Then*

$$R/I \cong \frac{R/J}{I/J}.$$

Theorem 16.34 Correspondence Theorem. *Let I be an ideal of a ring R. Then $S \mapsto S/I$ is a one-to-one correspondence between the set of subrings S containing I and the set of subrings of R/I. Furthermore, the ideals of R containing I correspond to ideals of R/I.*

16.4   Maximal and Prime Ideals

In this particular section we are especially interested in certain ideals of commutative rings. These ideals give us special types of factor rings. More specifically, we would like to characterize those ideals I of a commutative ring R such that R/I is an integral domain or a field.

A proper ideal M of a ring R is a **maximal ideal** of R if the ideal M is not a proper subset of any ideal of R except R itself. That is, M is a maximal ideal if for any ideal I properly containing M, $I = R$. The following theorem completely characterizes maximal ideals for commutative rings with identity in terms of their corresponding factor rings.

Theorem 16.35. *Let R be a commutative ring with identity and M an ideal in R. Then M is a maximal ideal of R if and only if R/M is a field.*

PROOF.   Let M be a maximal ideal in R. If R is a commutative ring, then R/M must also be a commutative ring. Clearly, $1 + M$ acts as an identity for R/M. We must also show that every nonzero element in R/M has an inverse. If $a + M$ is a nonzero element in R/M, then $a \notin M$. Define I to be the set $\{ra + m : r \in R \text{ and } m \in M\}$. We will show that I is an ideal in R. The set I is nonempty since $0a + 0 = 0$ is in I. If $r_1a + m_1$ and $r_2a + m_2$ are two elements in I, then

$$(r_1a + m_1) - (r_2a + m_2) = (r_1 - r_2)a + (m_1 - m_2)$$

is in I. Also, for any $r \in R$ it is true that $rI \subset I$; hence, I is closed under multiplication and satisfies the necessary conditions to be an ideal. Therefore, by Proposition 16.10 and the definition of an ideal, I is an ideal properly containing M. Since M is a maximal ideal, $I = R$; consequently, by the definition of I there must be an m in M and an element b in R such that $1 = ab + m$. Therefore,

$$1 + M = ab + M = ba + M = (a + M)(b + M).$$

Conversely, suppose that M is an ideal and R/M is a field. Since R/M is a field, it must contain at least two elements: $0 + M = M$ and $1 + M$. Hence, M is a proper ideal of R. Let I be any ideal properly containing M. We need to show that $I = R$. Choose a in I but not in M. Since $a + M$ is a nonzero element in a field, there exists an element $b + M$ in R/M such that $(a + M)(b + M) = ab + M = 1 + M$. Consequently, there exists an element $m \in M$ such that $ab + m = 1$ and 1 is in I. Therefore, $r1 = r \in I$ for all $r \in R$. Consequently, $I = R$. $\square$

Example 16.36. Let $p\mathbb{Z}$ be an ideal in $\mathbb{Z}$, where p is prime. Then $p\mathbb{Z}$ is a maximal ideal since $\mathbb{Z}/p\mathbb{Z} \cong \mathbb{Z}_p$ is a field.

A proper ideal P in a commutative ring R is called a ***prime ideal*** if whenever $ab \in P$, then either $a \in P$ or $b \in P$.[1]

Example 16.37. It is easy to check that the set $P = \{0, 2, 4, 6, 8, 10\}$ is an ideal in $\mathbb{Z}_{12}$. This ideal is prime. In fact, it is a maximal ideal.

Proposition 16.38. *Let R be a commutative ring with identity 1, where $1 \neq 0$. Then P is a prime ideal in R if and only if R/P is an integral domain.*

PROOF. First let us assume that P is an ideal in R and R/P is an integral domain. Suppose that $ab \in P$. If $a+P$ and $b+P$ are two elements of R/P such that $(a+P)(b+P) = 0+P = P$, then either $a + P = P$ or $b + P = P$. This means that either a is in P or b is in P, which shows that P must be prime.

Conversely, suppose that P is prime and

$$(a + P)(b + P) = ab + P = 0 + P = P.$$

Then $ab \in P$. If $a \notin P$, then b must be in P by the definition of a prime ideal; hence, $b + P = 0 + P$ and R/P is an integral domain. $\square$

Example 16.39. Every ideal in $\mathbb{Z}$ is of the form $n\mathbb{Z}$. The factor ring $\mathbb{Z}/n\mathbb{Z} \cong \mathbb{Z}_n$ is an integral domain only when n is prime. It is actually a field. Hence, the nonzero prime ideals in $\mathbb{Z}$ are the ideals $p\mathbb{Z}$, where p is prime. This example really justifies the use of the word "prime" in our definition of prime ideals.

Since every field is an integral domain, we have the following corollary.

Corollary 16.40. *Every maximal ideal in a commutative ring with identity is also a prime ideal.*

Historical Note

Amalie Emmy Noether, one of the outstanding mathematicians of the twentieth century, was born in Erlangen, Germany in 1882. She was the daughter of Max Noether (1844–1921), a distinguished mathematician at the University of Erlangen. Together with Paul Gordon (1837–1912), Emmy Noether's father strongly influenced her early education. She entered the University of Erlangen at the age of 18. Although women had been admitted to universities in England, France, and Italy for decades, there was great resistance to their presence at universities in Germany. Noether was one of only two women among the university's 986 students. After completing her doctorate under Gordon in 1907, she continued to do research at Erlangen, occasionally lecturing when her father was ill.

Noether went to Göttingen to study in 1916. David Hilbert and Felix Klein tried unsuccessfully to secure her an appointment at Göttingen. Some of the faculty objected to women lecturers, saying, "What will our soldiers think when they return to the university and are expected to learn at the feet of a woman?" Hilbert, annoyed at the question, responded, "Meine Herren, I do not see that the sex of a candidate is an argument against her admission as a Privatdozent. After all, the Senate is not a bathhouse." At the end of World War I, attitudes changed and conditions greatly improved for women. After Noether

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[1] It is possible to define prime ideals in a noncommutative ring. See [1] or [3].

passed her habilitation examination in 1919, she was given a title and was paid a small sum for her lectures.

In 1922, Noether became a Privatdozent at Göttingen. Over the next 11 years she used axiomatic methods to develop an abstract theory of rings and ideals. Though she was not good at lecturing, Noether was an inspiring teacher. One of her many students was B. L. van der Waerden, author of the first text treating abstract algebra from a modern point of view. Some of the other mathematicians Noether influenced or closely worked with were Alexandroff, Artin, Brauer, Courant, Hasse, Hopf, Pontryagin, von Neumann, and Weyl. One of the high points of her career was an invitation to address the International Congress of Mathematicians in Zurich in 1932. In spite of all the recognition she received from her colleagues, Noether's abilities were never recognized as they should have been during her lifetime. She was never promoted to full professor by the Prussian academic bureaucracy.

In 1933, Noether, a Jew, was banned from participation in all academic activities in Germany. She emigrated to the United States, took a position at Bryn Mawr College, and became a member of the Institute for Advanced Study at Princeton. Noether died suddenly on April 14, 1935. After her death she was eulogized by such notable scientists as Albert Einstein.

16.5   An Application to Software Design

The Chinese Remainder Theorem is a result from elementary number theory about the solution of systems of simultaneous congruences. The Chinese mathematician Sun-tsï wrote about the theorem in the first century A.D. This theorem has some interesting consequences in the design of software for parallel processors.

Lemma 16.41. *Let m and n be positive integers such that $\gcd(m, n) = 1$. Then for $a, b \in \mathbb{Z}$ the system*

$$x \equiv a \pmod{m}$$
$$x \equiv b \pmod{n}$$

has a solution. If x_1 and x_2 are two solutions of the system, then $x_1 \equiv x_2 \pmod{mn}$.

PROOF.   The equation $x \equiv a \pmod{m}$ has a solution since $a + km$ satisfies the equation for all $k \in \mathbb{Z}$. We must show that there exists an integer k_1 such that

$$a + k_1 m \equiv b \pmod{n}.$$

This is equivalent to showing that

$$k_1 m \equiv (b - a) \pmod{n}$$

has a solution for k_1. Since m and n are relatively prime, there exist integers s and t such that $ms + nt = 1$. Consequently,

$$(b - a)ms = (b - a) - (b - a)nt,$$

or

$$[(b - a)s]m \equiv (b - a) \pmod{n}.$$

Now let $k_1 = (b - a)s$.

To show that any two solutions are congruent modulo mn, let c_1 and c_2 be two solutions of the system. That is,

$$c_i \equiv a \pmod{m}$$

$$c_i \equiv b \pmod{n}$$

for $i = 1, 2$. Then

$$c_2 \equiv c_1 \pmod{m}$$
$$c_2 \equiv c_1 \pmod{n}.$$

Therefore, both m and n divide $c_1 - c_2$. Consequently, $c_2 \equiv c_1 \pmod{mn}$. □

Example 16.42. Let us solve the system

$$x \equiv 3 \pmod{4}$$
$$x \equiv 4 \pmod{5}.$$

Using the Euclidean algorithm, we can find integers s and t such that $4s + 5t = 1$. Two such integers are $s = 4$ and $t = -3$. Consequently,

$$x = a + k_1 m = 3 + 4k_1 = 3 + 4[(5 - 4)4] = 19.$$

Theorem 16.43 Chinese Remainder Theorem. *Let $n_1, n_2, \ldots, n_k$ be positive integers such that $\gcd(n_i, n_j) = 1$ for $i \neq j$. Then for any integers $a_1, \ldots, a_k$, the system*

$$x \equiv a_1 \pmod{n_1}$$
$$x \equiv a_2 \pmod{n_2}$$
$$\vdots$$
$$x \equiv a_k \pmod{n_k}$$

has a solution. Furthermore, any two solutions of the system are congruent modulo $n_1 n_2 \cdots n_k$.

PROOF. We will use mathematical induction on the number of equations in the system. If there are $k = 2$ equations, then the theorem is true by Lemma 16.41. Now suppose that the result is true for a system of k equations or less and that we wish to find a solution of

$$x \equiv a_1 \pmod{n_1}$$
$$x \equiv a_2 \pmod{n_2}$$
$$\vdots$$
$$x \equiv a_{k+1} \pmod{n_{k+1}}.$$

Considering the first k equations, there exists a solution that is unique modulo $n_1 \cdots n_k$, say a. Since $n_1 \cdots n_k$ and n_{k+1} are relatively prime, the system

$$x \equiv a \pmod{n_1 \cdots n_k}$$
$$x \equiv a_{k+1} \pmod{n_{k+1}}$$

has a solution that is unique modulo $n_1 \ldots n_{k+1}$ by the lemma. □

Example 16.44. Let us solve the system

$$x \equiv 3 \pmod{4}$$
$$x \equiv 4 \pmod{5}$$
$$x \equiv 1 \pmod{9}$$

$$x \equiv 5 \pmod 7.$$

From Example 16.42 we know that 19 is a solution of the first two congruences and any other solution of the system is congruent to 19 (mod 20). Hence, we can reduce the system to a system of three congruences:

$$x \equiv 19 \pmod{20}$$
$$x \equiv 1 \pmod 9$$
$$x \equiv 5 \pmod 7.$$

Solving the next two equations, we can reduce the system to

$$x \equiv 19 \pmod{180}$$
$$x \equiv 5 \pmod 7.$$

Solving this last system, we find that 19 is a solution for the system that is unique up to modulo 1260.

One interesting application of the Chinese Remainder Theorem in the design of computer software is that the theorem allows us to break up a calculation involving large integers into several less formidable calculations. A computer will handle integer calculations only up to a certain size due to the size of its processor chip, which is usually a 32 or 64-bit processor chip. For example, the largest integer available on a computer with a 64-bit processor chip is

$$2^{63} - 1 = 9{,}223{,}372{,}036{,}854{,}775{,}807.$$

Larger processors such as 128 or 256-bit have been proposed or are under development. There is even talk of a 512-bit processor chip. The largest integer that such a chip could store with be $2^{511} - 1$, which would be a 154 digit number. However, we would need to deal with much larger numbers to break sophisticated encryption schemes.

Special software is required for calculations involving larger integers which cannot be added directly by the machine. By using the Chinese Remainder Theorem we can break down large integer additions and multiplications into calculations that the computer can handle directly. This is especially useful on parallel processing computers which have the ability to run several programs concurrently.

Most computers have a single central processing unit (CPU) containing one processor chip and can only add two numbers at a time. To add a list of ten numbers, the CPU must do nine additions in sequence. However, a parallel processing computer has more than one CPU. A computer with 10 CPUs, for example, can perform 10 different additions at the same time. If we can take a large integer and break it down into parts, sending each part to a different CPU, then by performing several additions or multiplications simultaneously on those parts, we can work with an integer that the computer would not be able to handle as a whole.

Example 16.45. Suppose that we wish to multiply 2134 by 1531. We will use the integers 95, 97, 98, and 99 because they are relatively prime. We can break down each integer into four parts:

$$2134 \equiv 44 \pmod{95}$$
$$2134 \equiv 0 \pmod{97}$$
$$2134 \equiv 76 \pmod{98}$$
$$2134 \equiv 55 \pmod{99}$$

and

$$1531 \equiv 11 \pmod{95}$$
$$1531 \equiv 76 \pmod{97}$$
$$1531 \equiv 61 \pmod{98}$$
$$1531 \equiv 46 \pmod{99}.$$

Multiplying the corresponding equations, we obtain

$$2134 \cdot 1531 \equiv 44 \cdot 11 \equiv 9 \pmod{95}$$
$$2134 \cdot 1531 \equiv 0 \cdot 76 \equiv 0 \pmod{97}$$
$$2134 \cdot 1531 \equiv 76 \cdot 61 \equiv 30 \pmod{98}$$
$$2134 \cdot 1531 \equiv 55 \cdot 46 \equiv 55 \pmod{99}.$$

Each of these four computations can be sent to a different processor if our computer has several CPUs. By the above calculation, we know that $2134 \cdot 1531$ is a solution of the system

$$x \equiv 9 \pmod{95}$$
$$x \equiv 0 \pmod{97}$$
$$x \equiv 30 \pmod{98}$$
$$x \equiv 55 \pmod{99}.$$

The Chinese Remainder Theorem tells us that solutions are unique up to modulo $95 \cdot 97 \cdot 98 \cdot 99 = 89{,}403{,}930$. Solving this system of congruences for x tells us that $2134 \cdot 1531 = 3{,}267{,}154$.

The conversion of the computation into the four subcomputations will take some computing time. In addition, solving the system of congruences can also take considerable time. However, if we have many computations to be performed on a particular set of numbers, it makes sense to transform the problem as we have done above and to perform the necessary calculations simultaneously.

16.6 Exercises

1.   Which of the following sets are rings with respect to the usual operations of addition and multiplication? If the set is a ring, is it also a field?

(a) $7\mathbb{Z}$

(b) $\mathbb{Z}_{18}$

(c) $\mathbb{Q}(\sqrt{2}\,) = \{a + b\sqrt{2} : a, b \in \mathbb{Q}\}$

(d) $\mathbb{Q}(\sqrt{2}, \sqrt{3}\,) = \{a + b\sqrt{2} + c\sqrt{3} + d\sqrt{6} : a, b, c, d \in \mathbb{Q}\}$

(e) $\mathbb{Z}[\sqrt{3}\,] = \{a + b\sqrt{3} : a, b \in \mathbb{Z}\}$

(f) $R = \{a + b\sqrt[3]{3} : a, b \in \mathbb{Q}\}$

(g) $\mathbb{Z}[i] = \{a + bi : a, b \in \mathbb{Z} \text{ and } i^2 = -1\}$

(h) $\mathbb{Q}(\sqrt[3]{3}\,) = \{a + b\sqrt[3]{3} + c\sqrt[3]{9} : a, b, c \in \mathbb{Q}\}$

2.   Let R be the ring of 2×2 matrices of the form

$$\begin{pmatrix} a & b \\ 0 & 0 \end{pmatrix},$$

where $a, b \in \mathbb{R}$. Show that although R is a ring that has no identity, we can find a subring S of R with an identity.

3.   List or characterize all of the units in each of the following rings.

(a) $\mathbb{Z}_{10}$

(b) $\mathbb{Z}_{12}$

(c) $\mathbb{Z}_7$

(d) $\mathbb{M}_2(\mathbb{Z})$, the 2×2 matrices with entries in $\mathbb{Z}$

(e) $\mathbb{M}_2(\mathbb{Z}_2)$, the 2×2 matrices with entries in $\mathbb{Z}_2$

4.   Find all of the ideals in each of the following rings. Which of these ideals are maximal and which are prime?

(a) $\mathbb{Z}_{18}$

(b) $\mathbb{Z}_{25}$

(c) $\mathbb{M}_2(\mathbb{R})$, the 2×2 matrices with entries in $\mathbb{R}$

(d) $\mathbb{M}_2(\mathbb{Z})$, the 2×2 matrices with entries in $\mathbb{Z}$

(e) $\mathbb{Q}$

5.   For each of the following rings R with ideal I, give an addition table and a multiplication table for R/I.

(a) $R = \mathbb{Z}$ and $I = 6\mathbb{Z}$

(b) $R = \mathbb{Z}_{12}$ and $I = \{0, 3, 6, 9\}$

6.   Find all homomorphisms $\phi : \mathbb{Z}/6\mathbb{Z} \to \mathbb{Z}/15\mathbb{Z}$.

7.   Prove that $\mathbb{R}$ is not isomorphic to $\mathbb{C}$.

8.   Prove or disprove: The ring $\mathbb{Q}(\sqrt{2}) = \{a + b\sqrt{2} : a, b \in \mathbb{Q}\}$ is isomorphic to the ring $\mathbb{Q}(\sqrt{3}) = \{a + b\sqrt{3} : a, b \in \mathbb{Q}\}$.

9.   What is the characteristic of the field formed by the set of matrices

$$F = \left\{ \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}, \begin{pmatrix} 1 & 1 \\ 1 & 0 \end{pmatrix}, \begin{pmatrix} 0 & 1 \\ 1 & 1 \end{pmatrix}, \begin{pmatrix} 0 & 0 \\ 0 & 0 \end{pmatrix} \right\}$$

with entries in $\mathbb{Z}_2$?

10.   Define a map $\phi : \mathbb{C} \to \mathbb{M}_2(\mathbb{R})$ by

$$\phi(a + bi) = \begin{pmatrix} a & b \\ -b & a \end{pmatrix}.$$

Show that ϕ is an isomorphism of $\mathbb{C}$ with its image in $\mathbb{M}_2(\mathbb{R})$.

11.   Prove that the Gaussian integers, $\mathbb{Z}[i]$, are an integral domain.

12.   Prove that $\mathbb{Z}[\sqrt{3}\,i] = \{a + b\sqrt{3}\,i : a, b \in \mathbb{Z}\}$ is an integral domain.

13.   Solve each of the following systems of congruences.

(a)

$$x \equiv 2 \pmod 5$$
$$x \equiv 6 \pmod{11}$$

$$x \equiv 4 \pmod 7$$
$$x \equiv 7 \pmod 9$$
$$x \equiv 5 \pmod{11}$$

(b)

$$x \equiv 3 \pmod 7$$
$$x \equiv 0 \pmod 8$$
$$x \equiv 5 \pmod{15}$$

(d)

$$x \equiv 3 \pmod 5$$
$$x \equiv 0 \pmod 8$$
$$x \equiv 1 \pmod{11}$$
$$x \equiv 5 \pmod{13}$$

(c)

$$x \equiv 2 \pmod 4$$

14. Use the method of parallel computation outlined in the text to calculate $2234 + 4121$ by dividing the calculation into four separate additions modulo 95, 97, 98, and 99.

15. Explain why the method of parallel computation outlined in the text fails for $2134 \cdot 1531$ if we attempt to break the calculation down into two smaller calculations modulo 98 and 99.

16. If R is a field, show that the only two ideals of R are $\{0\}$ and R itself.

17. Let a be any element in a ring R with identity. Show that $(-1)a = -a$.

18. Let $\phi : R \to S$ be a ring homomorphism. Prove each of the following statements.

(a) If R is a commutative ring, then $\phi(R)$ is a commutative ring.

(b) $\phi(0) = 0$.

(c) Let 1_R and 1_S be the identities for R and S, respectively. If ϕ is onto, then $\phi(1_R) = 1_S$.

(d) If R is a field and $\phi(R) \neq 0$, then $\phi(R)$ is a field.

19. Prove that the associative law for multiplication and the distributive laws hold in R/I.

20. Prove the Second Isomorphism Theorem for rings: Let I be a subring of a ring R and J an ideal in R. Then $I \cap J$ is an ideal in I and

$$I/I \cap J \cong I + J/J.$$

21. Prove the Third Isomorphism Theorem for rings: Let R be a ring and I and J be ideals of R, where $J \subset I$. Then

$$R/I \cong \frac{R/J}{I/J}.$$

22. Prove the Correspondence Theorem: Let I be an ideal of a ring R. Then $S \to S/I$ is a one-to-one correspondence between the set of subrings S containing I and the set of subrings of R/I. Furthermore, the ideals of R correspond to ideals of R/I.

23. Let R be a ring and S a subset of R. Show that S is a subring of R if and only if each of the following conditions is satisfied.

(a) $S \neq \emptyset$.

(b) $rs \in S$ for all $r, s \in S$.

(c) $r - s \in S$ for all $r, s \in S$.

24.   Let R be a ring with a collection of subrings $\{R_\alpha\}$. Prove that $\bigcap R_\alpha$ is a subring of R. Give an example to show that the union of two subrings is not necessarily a subring.

25.   Let $\{I_\alpha\}_{\alpha \in A}$ be a collection of ideals in a ring R. Prove that $\bigcap_{\alpha \in A} I_\alpha$ is also an ideal in R. Give an example to show that if I_1 and I_2 are ideals in R, then $I_1 \cup I_2$ may not be an ideal.

26.   Let R be an integral domain. Show that if the only ideals in R are $\{0\}$ and R itself, R must be a field.

27.   Let R be a commutative ring. An element a in R is **nilpotent** if $a^n = 0$ for some positive integer n. Show that the set of all nilpotent elements forms an ideal in R.

28.   A ring R is a **Boolean ring** if for every $a \in R$, $a^2 = a$. Show that every Boolean ring is a commutative ring.

29.   Let R be a ring, where $a^3 = a$ for all $a \in R$. Prove that R must be a commutative ring.

30.   Let R be a ring with identity 1_R and S a subring of R with identity 1_S. Prove or disprove that $1_R = 1_S$.

31.   If we do not require the identity of a ring to be distinct from 0, we will not have a very interesting mathematical structure. Let R be a ring such that $1 = 0$. Prove that $R = \{0\}$.

32.   Let S be a nonempty subset of a ring R. Prove that there is a subring R' of R that contains S.

33.   Let R be a ring. Define the **center** of R to be

$$Z(R) = \{a \in R : ar = ra \text{ for all } r \in R\}.$$

Prove that $Z(R)$ is a commutative subring of R.

34.   Let p be prime. Prove that

$$\mathbb{Z}_{(p)} = \{a/b : a, b \in \mathbb{Z} \text{ and } \gcd(b, p) = 1\}$$

is a ring. The ring $\mathbb{Z}_{(p)}$ is called the **ring of integers localized at** p.

35.   Prove or disprove: Every finite integral domain is isomorphic to $\mathbb{Z}_p$.

36.   Let R be a ring with identity.

(a) Let u be a unit in R. Define a map $i_u : R \to R$ by $r \mapsto uru^{-1}$. Prove that i_u is an automorphism of R. Such an automorphism of R is called an inner automorphism of R. Denote the set of all inner automorphisms of R by $\mathrm{Inn}(R)$.

(b) Denote the set of all automorphisms of R by $\mathrm{Aut}(R)$. Prove that $\mathrm{Inn}(R)$ is a normal subgroup of $\mathrm{Aut}(R)$.

(c) Let $U(R)$ be the group of units in R. Prove that the map

$$\phi : U(R) \to \mathrm{Inn}(R)$$

defined by $u \mapsto i_u$ is a homomorphism. Determine the kernel of ϕ.

(d) Compute $\mathrm{Aut}(\mathbb{Z})$, $\mathrm{Inn}(\mathbb{Z})$, and $U(\mathbb{Z})$.

37.   Let R and S be arbitrary rings. Show that their Cartesian product is a ring if we define addition and multiplication in $R \times S$ by

(a) $(r, s) + (r', s') = (r + r', s + s')$

(b) $(r, s)(r', s') = (rr', ss')$

38.   An element x in a ring is called an ***idempotent*** if $x^2 = x$. Prove that the only idempotents in an integral domain are 0 and 1. Find a ring with a idempotent x not equal to 0 or 1.

39.   Let $\gcd(a, n) = d$ and $\gcd(b, d) \neq 1$. Prove that $ax \equiv b \pmod{n}$ does not have a solution.

40. **The Chinese Remainder Theorem for Rings.**   Let R be a ring and I and J be ideals in R such that $I + J = R$.

(a) Show that for any r and s in R, the system of equations

$$x \equiv r \pmod{I}$$
$$x \equiv s \pmod{J}$$

has a solution.

(b) In addition, prove that any two solutions of the system are congruent modulo $I \cap J$.

(c) Let I and J be ideals in a ring R such that $I + J = R$. Show that there exists a ring isomorphism
$$R/(I \cap J) \cong R/I \times R/J.$$

16.7   Programming Exercise

1.   Write a computer program implementing fast addition and multiplication using the Chinese Remainder Theorem and the method outlined in the text.

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16.8   Sage

Rings are very important in your study of abstract algebra, and similarly, they are very
important in the design and use of Sage. There is a lot of material in this chapter, and
there are many corresponding commands in Sage.

Creating Rings

Here is a list of various rings, domains and fields you can construct simply.

1. `Integers()`, `ZZ`: the integral domain of positive and negative integers, $\mathbb{Z}$.

2. `Integers(n)`: the integers mod n, $\mathbb{Z}_n$. A field when n is prime, but just a ring for
    composite n.

3. `QQ`: the field of rational numbers, $\mathbb{Q}$.

4. `RR`, `CC`: the field of real numbers and the field of complex numbers, $\mathbb{R}$, $\mathbb{C}$. It is impossible
    to create *every* real number inside a computer, so technically these sets do not behave
    as fields, but only give a good imitiation of the real thing. We say they are ***inexact***
    rings to make this point.

5. `QuadraticField(n)`: the field formed by combining the rationals with a solution to
    the polynomial equation $x^2 - n = 0$. The notation in the text is $\mathbb{Q}[\sqrt{n}]$. A functional
    equivalent can be made with the syntax `QQ[sqrt(n)]`. Note that `n` can be negative.

6. `CyclotomicField(n)`: the field formed by combining the rationals with the solutions
    to the polynomial equation $x^n - 1 = 0$.

7. `QQbar`: the field formed by combining the rationals with the solutions to *every* poly-
    nomial equation with integer coefficients. This is known as a the field of algebraic
    numbers, denoted as $\overline{\mathbb{Q}}$.

8. `FiniteField(p)`: for a prime p, the field of integers $\mathbb{Z}_p$.

  If you print a description of some of the above rings, you will sometimes see a new
symbol introduced. Consider the following example:

```
F = QuadraticField(7)
F
```

```
Number Field in a with defining polynomial x^2 - 7
```

```
root = F.gen(0)
root^2
```

```
root
```

a

```
(2*root)^3
```

56*a

Here `Number Field` describes an object generally formed by combining the rationals with another number (here $\sqrt{7}$). "a" is a new symbol which behaves as a root of the polynomial $x^2 - 7$. We do not say which root, $\sqrt{7}$ or $-\sqrt{7}$, and as we understand the theory better we will see that this does not really matter.

We can obtain this root as a generator of the number field, and then manipulate it. First squaring `root` yields 7. Notice that `root` prints as `a`. Notice, too, that computations with `root` behave as if it was *either* root of $x^2 - 7$, and results print using `a`.

This can get a bit confusing, inputing computations with `root` and getting output in terms of `a`. Fortunately, there is a better way. Consider the following example:

```
F.<b> = QuadraticField(7)
F
```

Number Field **in** b with defining polynomial x^2 - 7

```
b^2
```

7

```
(2*b)^3
```

56*b

With the syntax `F.<b>` we can create the field `F` along with specifying a generator `b` using a name of our choosing. Then computations can use `b` in both input and output as a root of $x^2 - 7$.

Here are three new rings that are best created using this new syntax.

1. `F.<a> = FiniteField(p^n)`: We will later have a theorem that tells us that finite fields only exist with orders equal to to a power of a prime. When the power is larger than 1, then we need a generator, here given as `a`.

2. `P.<x>=R[]`: the ring of all polynomials in the variable `x`, with coefficients from the ring `R`. Notice that `R` can be *any* ring, so this is a very general construction that uses one ring to form another. See an example below.

3. `Q.<r,s,t> = QuaternionAlgebra(n, m)`: the rationals combined with indeterminates `r`, `s` and `t` such that $r^2 = n$, $s^2 = m$ and $t = rs = -sr$. This is a generalization of the quaternions described in this chapter, though over the rationals rather than the reals, so it is an exact ring. Notice that this is one of the few noncommutative rings in Sage. The "usual" quaternions would be constructed with `Q.<I,J,K> = QuaternionAlgebra(-1, -1)`. (Notice that using `I` here is not a good choice, because it will then clobber the symbol `I` used for complex numbers.)

Syntax specifying names for generators can be used for many of the above rings as well, such as demonstrated above for quadratic fields and below for cyclotomic fields.

```
C.<t> = CyclotomicField(8)
C.random_element()
```

```
-2/11*t^2 + t - 1
```

Properties of Rings

The examples below demonstrate how to query certain properties of rings. If you are playing along, be sure to execute the first compute cell to define the various rings involved in the examples.

```
Z7 = Integers(7)
Z9 = Integers(9)
Q = QuadraticField(-11)
F.<a> = FiniteField(3^2)
P.<x> = Z7[]
S.<f,g,h> = QuaternionAlgebra(-7, 3)
```

Exact versus inexact.

```
QQ.is_exact()
```

```
True
```

```
RR.is_exact()
```

```
False
```

Finite versus infinite.

```
Z7.is_finite()
```

```
True
```

```
Z7.is_finite()
```

```
True
```

Integral domain?

```
Z7.is_integral_domain()
```

```
True
```

```
Z9.is_integral_domain()
```

```
False
```

Field?

```
Z9.is_field()
```

```
False
```

```
F.is_field()
```

```
True
```

```
Q.is_field()
```

True

Commutative?

```
Q.is_commutative()
```

True

```
S.is_commutative()
```

False

Characteristic.

```
Z7.characteristic()
```

7

```
Z9.characteristic()
```

9

```
Q.characteristic()
```

0

```
F.characteristic()
```

3

```
P.characteristic()
```

7

```
S.characteristic()
```

0

Additive and multiplicative identities *print* like you would expect, but notice that while they may *print* identically, they could be *different* because of the ring they live in.

```
b = Z9.zero(); b
```

0

```
b.parent()
```

Ring of integers modulo 9

```
c = Q.zero(); c
```

0

```
c.parent()
```

Number Field **in** a with defining polynomial x^2 + 11

```
b == c
```

False

```
d = Z9.one(); d
```

1

```
d.parent()
```

Ring of integers modulo 9

```
e = Q.one(); e
```

1

```
e.parent()
```

Number Field **in** a with defining polynomial x^2 + 11

```
d == e
```

False

There is some support for subrings. For example, Q and S are extensions of the rationals, while F is totally distinct from the rationals.

```
QQ.is_subring(Q)
```

True

```
QQ.is_subring(S)
```

True

```
QQ.is_subring(F)
```

False

Not every element of a ring may have a multiplicative inverse, in other words, not every element has to be a unit (unless the ring is a field). It would now be good practice to check if an element is a unit before you try computing its inverse.

```
three = Z9(3)
three.is_unit()
```

False

```
three*three
```

0

```
four = Z9(4)
four.is_unit()
```

True

```
g = four^-1; g
```

7

```
four*g
```

1

Quotient Structure

Ideals are the normal subgroups of rings and allow us to build "quotients" — basically new rings defined on equivalence classes of elements of the original ring. Sage support for ideals is variable. When they can be created, there is not always a lot you can do with them. But they work well in certain very important cases.

The integers, $\mathbb{Z}$, have ideals that are just multiples of a single integer. We can create them with the `.ideal()` method or just by wrting a scalar multiple of `ZZ`. And then the quotient is isomorphic to a well-understood ring. (Notice that `I` is a bad name for an ideal if we want to work with complex numbers later.)

```
I1 = ZZ.ideal(4)
I2 = 4*ZZ
I3 = (-4)*ZZ
I1 == I2
```

True

```
I2 == I3
```

True

```
Q = ZZ.quotient(I1); Q
```

Ring of integers modulo 4

```
Q == Integers(4)
```

True

We might normally be more careful about the last statement. The quotient is a set of equivalence classes, each infinite, and certainly not a single integer. But the quotient is *isomorphic* to $\mathbb{Z}_4$, so Sage just makes this identification.

```
Z7 = Integers(7)
P.<y> = Z7[]
M = P.ideal(y^2+4)
Q = P.quotient(M)
Q
```

Univariate Quotient Polynomial Ring **in** ybar over
Ring of integers modulo 7 with modulus y^2 + 4

```
Q.random_element()
```

2*ybar + 6

```
Q.order()
```

49

```
Q.is_field()
```

True

Notice that the construction of the quotient ring has created a new generator, converting y (y) to ybar ($\overline{y}$). We can override this as before with the syntax demonstrated below.

```
Q.<t> = P.quotient(M); Q
```

```
Univariate Quotient Polynomial Ring in t over
Ring of integers modulo 7 with modulus y^2 + 4
```

```
Q.random_element()
```

4*t + 6

So from a quotient of an infinite ring and an ideal (which is also a ring), we create a field, which is finite. Understanding this construction will be an important theme in the next few chapters. To see how remarkable it is, consider what happens with just one little change.

```
Z7 = Integers(7)
P.<y> = Z7[]
M = P.ideal(y^2+3)
Q.<t> = P.quotient(M)
Q
```

```
Univariate Quotient Polynomial Ring in t over
Ring of integers modulo 7 with modulus y^2 + 3
```

```
Q.random_element()
```

3*t + 1

```
Q.order()
```

49

```
Q.is_field()
```

False

There are a few methods available which will give us properties of ideals. In particular, we can check for prime and maximal ideals in rings of polynomials. Examine the results above and below in the context of Theorem 16.35.

```
Z7 = Integers(7)
P.<y> = Z7[]
M = P.ideal(y^2+4)
N = P.ideal(y^2+3)
M.is_maximal()
```

True

```
N.is_maximal()
```

```
False
```

The fact that `M` is a prime ideal is verification of Corollary 16.40.

```
M.is_prime()
```

```
True
```

```
N.is_prime()
```

```
False
```

Ring Homomorphisms

When Sage is presented with `3 + 4/3`, how does it know that `3` is meant to be an integer? And then to add it to a rational, how does it know that we really want to view the computation as $3/1 + 4/3$? This is really easy for you and me, but devilishly hard for a program, and you can imagine it getting ever more complicated with the many possible rings in Sage, subrings, matrices, etc. Part of the answer is that Sage uses ring homomorphisms to "translate" objects (numbers) between rings.

We will give an example below, but not pursue the topic much further. For the curious, reading the Sage documentation and experimenting would be a good exercise.

```
H = Hom(ZZ, QQ)
phi = H([1])
phi
```

```
Ring morphism:
  From: Integer Ring
  To:   Rational Field
  Defn: 1 |--> 1
```

```
phi.parent()
```

```
Set of Homomorphisms from Integer Ring to Rational Field
```

```
a = 3; a
```

```
3
```

```
a.parent()
```

```
Integer Ring
```

```
b = phi(3); b
```

```
3
```

```
b.parent()
```

```
Rational Field
```

So `phi` is a homomorphism ("morphism") that converts integers (the domain is `ZZ`) into rationals (the codomain is `QQ`), whose parent is a set of homomorphisms that Sage calls a "homset." Even though `a` and `b` both print as `3`, which is indistinguishable to our eyes, the parents of `a` and `b` are different. Yet the numerical value of the two objects has not changed.

16.9   Sage Exercises

1.   Define the two rings $\mathbb{Z}_{11}$ and $\mathbb{Z}_{12}$ with the commands `R = Integers(11)` and `S = Integers(12)`. For each ring, use the relevant command to determine: if the ring is finite, if it is commutative, if it is an integral domain and if it is a field. Then use single Sage commands to find the order of the ring, list the elements, and output the multiplicative identity (i.e. 1, if it exists).

2.   Define `R` to be the ring of integers, $\mathbb{Z}$, by executing `R = ZZ` or `R = Integers()`. A command like `R.ideal(4)` will create the principal ideal $\langle 4 \rangle$. The same command can accept more than one generator, so for example, `R.ideal(3, 5)` will create the ideal $\{a \cdot 3 + b \cdot 5 \mid a, b \in \mathbb{Z}\}$. Create several ideals of $\mathbb{Z}$ with two generators and ask Sage to print each as you create it. Explain what you observe and then create code that will test your observation for thousands of different examples.

3.   Create a finite field F of order 81 with `F.<t>=FiniteField(3^4)`.

(a) List the elements of F.

(b) Obtain the generators of F with `F.gens()`.

(c) Obtain the first generator of F and save it as `u` with `u = F.0` (alternatively, `u = F.gen(0)`).

(d) Compute the first 80 powers of `u` and comment.

(e) The generator you have worked with above is a root of a polynomial over $\mathbb{Z}_3$. Obtain this polynomial with `F.modulus()` and use this observation to explain the entry in your list of powers that is the fourth power of the generator.

4.   Build and analyze a quotient ring as follows:

(a) Use `P.<z>=Integers(7)[]` to construct a ring P of polynomials in z with coefficients from $\mathbb{Z}_7$.

(b) Use `K = P.ideal(z^2+z+3)` to build a principal ideal K generated by the polynomial $z^2 + z + 3$.

(c) Use `H = P.quotient(K)` to build H, the quotient ring of P by K.

(d) Use Sage to verify that H is a field.

(e) As in the previous exercise, obtain a generator and examine the proper collection of powers of that generator.

$$\boxed{17}$$

Polynomials

Most people are fairly familiar with polynomials by the time they begin to study abstract algebra. When we examine polynomial expressions such as

$$p(x) = x^3 - 3x + 2$$
$$q(x) = 3x^2 - 6x + 5,$$

we have a pretty good idea of what $p(x) + q(x)$ and $p(x)q(x)$ mean. We just add and multiply polynomials as functions; that is,

$$\begin{aligned}(p+q)(x) &= p(x) + q(x) \\ &= (x^3 - 3x + 2) + (3x^2 - 6x + 5) \\ &= x^3 + 3x^2 - 9x + 7\end{aligned}$$

and

$$\begin{aligned}(pq)(x) &= p(x)q(x) \\ &= (x^3 - 3x + 2)(3x^2 - 6x + 5) \\ &= 3x^5 - 6x^4 - 4x^3 + 24x^2 - 27x + 10.\end{aligned}$$

It is probably no surprise that polynomials form a ring. In this chapter we shall emphasize the algebraic structure of polynomials by studying polynomial rings. We can prove many results for polynomial rings that are similar to the theorems we proved for the integers. Analogs of prime numbers, the division algorithm, and the Euclidean algorithm exist for polynomials.

17.1 Polynomial Rings

Throughout this chapter we shall assume that R is a commutative ring with identity. Any expression of the form

$$f(x) = \sum_{i=0}^{n} a_i x^i = a_0 + a_1 x + a_2 x^2 + \cdots + a_n x^n,$$

where $a_i \in R$ and $a_n \neq 0$, is called a **polynomial over** R with **indeterminate** x. The elements $a_0, a_1, \ldots, a_n$ are called the **coefficients** of f. The coefficient a_n is called the **leading coefficient**. A polynomial is called **monic** if the leading coefficient is 1. If n is the largest nonnegative number for which $a_n \neq 0$, we say that the **degree** of f is n and write $\deg f(x) = n$. If no such n exists—that is, if $f = 0$ is the zero polynomial—then the degree

of f is defined to be $-\infty$. We will denote the set of all polynomials with coefficients in a ring R by $R[x]$. Two polynomials are equal exactly when their corresponding coefficients are equal; that is, if we let

$$p(x) = a_0 + a_1 x + \cdots + a_n x^n$$
$$q(x) = b_0 + b_1 x + \cdots + b_m x^m,$$

then $p(x) = q(x)$ if and only if $a_i = b_i$ for all $i \geq 0$.

To show that the set of all polynomials forms a ring, we must first define addition and multiplication. We define the sum of two polynomials as follows. Let

$$p(x) = a_0 + a_1 x + \cdots + a_n x^n$$
$$q(x) = b_0 + b_1 x + \cdots + b_m x^m.$$

Then the sum of $p(x)$ and $q(x)$ is

$$p(x) + q(x) = c_0 + c_1 x + \cdots + c_k x^k,$$

where $c_i = a_i + b_i$ for each i. We define the product of $p(x)$ and $q(x)$ to be

$$p(x)q(x) = c_0 + c_1 x + \cdots + c_{m+n} x^{m+n},$$

where

$$c_i = \sum_{k=0}^{i} a_k b_{i-k} = a_0 b_i + a_1 b_{i-1} + \cdots + a_{i-1} b_1 + a_i b_0$$

for each i. Notice that in each case some of the coefficients may be zero.

Example 17.1. Suppose that

$$p(x) = 3 + 0x + 0x^2 + 2x^3 + 0x^4$$

and

$$q(x) = 2 + 0x - x^2 + 0x^3 + 4x^4$$

are polynomials in $\mathbb{Z}[x]$. If the coefficient of some term in a polynomial is zero, then we usually just omit that term. In this case we would write $p(x) = 3 + 2x^3$ and $q(x) = 2 - x^2 + 4x^4$. The sum of these two polynomials is

$$p(x) + q(x) = 5 - x^2 + 2x^3 + 4x^4.$$

The product,

$$p(x)q(x) = (3 + 2x^3)(2 - x^2 + 4x^4) = 6 - 3x^2 + 4x^3 + 12x^4 - 2x^5 + 8x^7,$$

can be calculated either by determining the c_is in the definition or by simply multiplying polynomials in the same way as we have always done.

Example 17.2. Let

$$p(x) = 3 + 3x^3 \qquad \text{and} \qquad q(x) = 4 + 4x^2 + 4x^4$$

be polynomials in $\mathbb{Z}_{12}[x]$. The sum of $p(x)$ and $q(x)$ is $7 + 4x^2 + 3x^3 + 4x^4$. The product of the two polynomials is the zero polynomial. This example tells us that we can not expect $R[x]$ to be an integral domain if R is not an integral domain.

Theorem 17.3. *Let R be a commutative ring with identity. Then $R[x]$ is a commutative ring with identity.*

PROOF.   Our first task is to show that $R[x]$ is an abelian group under polynomial addition. The zero polynomial, $f(x) = 0$, is the additive identity. Given a polynomial $p(x) = \sum_{i=0}^{n} a_i x^i$, the inverse of $p(x)$ is easily verified to be $-p(x) = \sum_{i=0}^{n} (-a_i) x^i = -\sum_{i=0}^{n} a_i x^i$. Commutativity and associativity follow immediately from the definition of polynomial addition and from the fact that addition in R is both commutative and associative.

To show that polynomial multiplication is associative, let

$$p(x) = \sum_{i=0}^{m} a_i x^i,$$

$$q(x) = \sum_{i=0}^{n} b_i x^i,$$

$$r(x) = \sum_{i=0}^{p} c_i x^i.$$

Then

$$
\begin{aligned}
[p(x)q(x)]r(x) &= \left[\left(\sum_{i=0}^{m} a_i x^i \right) \left(\sum_{i=0}^{n} b_i x^i \right) \right] \left(\sum_{i=0}^{p} c_i x^i \right) \\
&= \left[\sum_{i=0}^{m+n} \left(\sum_{j=0}^{i} a_j b_{i-j} \right) x^i \right] \left(\sum_{i=0}^{p} c_i x^i \right) \\
&= \sum_{i=0}^{m+n+p} \left[\sum_{j=0}^{i} \left(\sum_{k=0}^{j} a_k b_{j-k} \right) c_{i-j} \right] x^i \\
&= \sum_{i=0}^{m+n+p} \left(\sum_{j+k+l=i} a_j b_k c_l \right) x^i \\
&= \sum_{i=0}^{m+n+p} \left[\sum_{j=0}^{i} a_j \left(\sum_{k=0}^{i-j} b_k c_{i-j-k} \right) \right] x^i \\
&= \left(\sum_{i=0}^{m} a_i x^i \right) \left[\sum_{i=0}^{n+p} \left(\sum_{j=0}^{i} b_j c_{i-j} \right) x^i \right] \\
&= \left(\sum_{i=0}^{m} a_i x^i \right) \left[\left(\sum_{i=0}^{n} b_i x^i \right) \left(\sum_{i=0}^{p} c_i x^i \right) \right] \\
&= p(x)[q(x)r(x)]
\end{aligned}
$$

The commutativity and distribution properties of polynomial multiplication are proved in a similar manner. We shall leave the proofs of these properties as an exercise.   □

Proposition 17.4. *Let $p(x)$ and $q(x)$ be polynomials in $R[x]$, where R is an integral domain. Then $\deg p(x) + \deg q(x) = \deg(p(x)q(x))$. Furthermore, $R[x]$ is an integral domain.*

PROOF.   Suppose that we have two nonzero polynomials

$$p(x) = a_m x^m + \cdots + a_1 x + a_0$$

and
$$q(x) = b_n x^n + \cdots + b_1 x + b_0$$

with $a_m \neq 0$ and $b_n \neq 0$. The degrees of $p(x)$ and $q(x)$ are m and n, respectively. The leading term of $p(x)q(x)$ is $a_m b_n x^{m+n}$, which cannot be zero since R is an integral domain; hence, the degree of $p(x)q(x)$ is $m + n$, and $p(x)q(x) \neq 0$. Since $p(x) \neq 0$ and $q(x) \neq 0$ imply that $p(x)q(x) \neq 0$, we know that $R[x]$ must also be an integral domain. $\square$

We also want to consider polynomials in two or more variables, such as $x^2 - 3xy + 2y^3$. Let R be a ring and suppose that we are given two indeterminates x and y. Certainly we can form the ring $(R[x])[y]$. It is straightforward but perhaps tedious to show that $(R[x])[y] \cong R([y])[x]$. We shall identify these two rings by this isomorphism and simply write $R[x, y]$. The ring $R[x, y]$ is called the ***ring of polynomials in two indeterminates x and y with coefficients in*** R. We can define the ***ring of polynomials in*** n ***indeterminates with coefficients in*** R similarly. We shall denote this ring by $R[x_1, x_2, \ldots, x_n]$.

Theorem 17.5. *Let R be a commutative ring with identity and $\alpha \in R$. Then we have a ring homomorphism $\phi_\alpha : R[x] \to R$ defined by*

$$\phi_\alpha(p(x)) = p(\alpha) = a_n \alpha^n + \cdots + a_1 \alpha + a_0,$$

where $p(x) = a_n x^n + \cdots + a_1 x + a_0$.

PROOF.   Let $p(x) = \sum_{i=0}^{n} a_i x^i$ and $q(x) = \sum_{i=0}^{m} b_i x^i$. It is easy to show that $\phi_\alpha(p(x) + q(x)) = \phi_\alpha(p(x)) + \phi_\alpha(q(x))$. To show that multiplication is preserved under the map ϕ_α, observe that

$$
\begin{aligned}
\phi_\alpha(p(x))\phi_\alpha(q(x)) &= p(\alpha)q(\alpha) \\
&= \left(\sum_{i=0}^{n} a_i \alpha^i \right) \left(\sum_{i=0}^{m} b_i \alpha^i \right) \\
&= \sum_{i=0}^{m+n} \left(\sum_{k=0}^{i} a_k b_{i-k} \right) \alpha^i \\
&= \phi_\alpha(p(x)q(x)).
\end{aligned}
$$
$\square$

The map $\phi_\alpha : R[x] \to R$ is called the ***evaluation homomorphism*** at α.

17.2   The Division Algorithm

Recall that the division algorithm for integers (Theorem 2.9) says that if a and b are integers with $b > 0$, then there exist unique integers q and r such that $a = bq + r$, where $0 \leq r < b$. The algorithm by which q and r are found is just long division. A similar theorem exists for polynomials. The division algorithm for polynomials has several important consequences. Since its proof is very similar to the corresponding proof for integers, it is worthwhile to review Theorem 2.9 at this point.

Theorem 17.6 Division Algorithm. *Let $f(x)$ and $g(x)$ be polynomials in $F[x]$, where F is a field and $g(x)$ is a nonzero polynomial. Then there exist unique polynomials $q(x), r(x) \in F[x]$ such that*

$$f(x) = g(x)q(x) + r(x),$$

where either $\deg r(x) < \deg g(x)$ or $r(x)$ is the zero polynomial.

PROOF.   We will first consider the existence of $q(x)$ and $r(x)$. If $f(x)$ is the zero polynomial, then

$$0 = 0 \cdot g(x) + 0;$$

hence, both q and r must also be the zero polynomial. Now suppose that $f(x)$ is not the zero polynomial and that $\deg f(x) = n$ and $\deg g(x) = m$. If $m > n$, then we can let $q(x) = 0$ and $r(x) = f(x)$. Hence, we may assume that $m \leq n$ and proceed by induction on n. If

$$f(x) = a_n x^n + a_{n-1} x^{n-1} + \cdots + a_1 x + a_0$$
$$g(x) = b_m x^m + b_{m-1} x^{m-1} + \cdots + b_1 x + b_0$$

the polynomial

$$f'(x) = f(x) - \frac{a_n}{b_m} x^{n-m} g(x)$$

has degree less than n or is the zero polynomial. By induction, there exist polynomials $q'(x)$ and $r(x)$ such that

$$f'(x) = q'(x) g(x) + r(x),$$

where $r(x) = 0$ or the degree of $r(x)$ is less than the degree of $g(x)$. Now let

$$q(x) = q'(x) + \frac{a_n}{b_m} x^{n-m}.$$

Then

$$f(x) = g(x) q(x) + r(x),$$

with $r(x)$ the zero polynomial or $\deg r(x) < \deg g(x)$.

To show that $q(x)$ and $r(x)$ are unique, suppose that there exist two other polynomials $q_1(x)$ and $r_1(x)$ such that $f(x) = g(x) q_1(x) + r_1(x)$ with $\deg r_1(x) < \deg g(x)$ or $r_1(x) = 0$, so that

$$f(x) = g(x) q(x) + r(x) = g(x) q_1(x) + r_1(x),$$

and

$$g(x)[q(x) - q_1(x)] = r_1(x) - r(x).$$

If $q(x) - q_1(x)$ is not the zero polynomial, then

$$\deg(g(x)[q(x) - q_1(x)]) = \deg(r_1(x) - r(x)) \geq \deg g(x).$$

However, the degrees of both $r(x)$ and $r_1(x)$ are strictly less than the degree of $g(x)$; therefore, $r(x) = r_1(x)$ and $q(x) = q_1(x)$. $\square$

Example 17.7. The division algorithm merely formalizes long division of polynomials, a task we have been familiar with since high school. For example, suppose that we divide $x^3 - x^2 + 2x - 3$ by $x - 2$.

$$
\begin{array}{r}
x^2 + x + 4 \\
x - 2 \overline{\smash{\big)}\ x^3 - x^2 + 2x - 3} \\
\underline{x^3 - 2x^2} \\
x^2 + 2x - 3 \\
\underline{x^2 - 2x} \\
4x - 3 \\
\underline{4x - 8} \\
5
\end{array}
$$

Hence, $x^3 - x^2 + 2x - 3 = (x - 2)(x^2 + x + 4) + 5$.

Let $p(x)$ be a polynomial in $F[x]$ and $\alpha \in F$. We say that α is a **zero** or **root** of $p(x)$ if $p(x)$ is in the kernel of the evaluation homomorphism ϕ_α. All we are really saying here is that α is a zero of $p(x)$ if $p(\alpha) = 0$.

Corollary 17.8. *Let F be a field. An element $\alpha \in F$ is a zero of $p(x) \in F[x]$ if and only if $x - \alpha$ is a factor of $p(x)$ in $F[x]$.*

PROOF.   Suppose that $\alpha \in F$ and $p(\alpha) = 0$. By the division algorithm, there exist polynomials $q(x)$ and $r(x)$ such that

$$p(x) = (x - \alpha)q(x) + r(x)$$

and the degree of $r(x)$ must be less than the degree of $x - \alpha$. Since the degree of $r(x)$ is less than 1, $r(x) = a$ for $a \in F$; therefore,

$$p(x) = (x - \alpha)q(x) + a.$$

But

$$0 = p(\alpha) = 0 \cdot q(\alpha) + a = a;$$

consequently, $p(x) = (x - \alpha)q(x)$, and $x - \alpha$ is a factor of $p(x)$.

Conversely, suppose that $x - \alpha$ is a factor of $p(x)$; say $p(x) = (x - \alpha)q(x)$. Then $p(\alpha) = 0 \cdot q(\alpha) = 0$. $\square$

Corollary 17.9. *Let F be a field. A nonzero polynomial $p(x)$ of degree n in $F[x]$ can have at most n distinct zeros in F.*

PROOF.   We will use induction on the degree of $p(x)$. If $\deg p(x) = 0$, then $p(x)$ is a constant polynomial and has no zeros. Let $\deg p(x) = 1$. Then $p(x) = ax + b$ for some a and b in F. If α_1 and α_2 are zeros of $p(x)$, then $a\alpha_1 + b = a\alpha_2 + b$ or $\alpha_1 = \alpha_2$.

Now assume that $\deg p(x) > 1$. If $p(x)$ does not have a zero in F, then we are done. On the other hand, if α is a zero of $p(x)$, then $p(x) = (x - \alpha)q(x)$ for some $q(x) \in F[x]$ by Corollary 17.8. The degree of $q(x)$ is $n - 1$ by Proposition 17.4. Let β be some other zero of $p(x)$ that is distinct from α. Then $p(\beta) = (\beta - \alpha)q(\beta) = 0$. Since $\alpha \neq \beta$ and F is a field, $q(\beta) = 0$. By our induction hypothesis, $q(x)$ can have at most $n - 1$ zeros in F that are distinct from α. Therefore, $p(x)$ has at most n distinct zeros in F. $\square$

Let F be a field. A monic polynomial $d(x)$ is a **greatest common divisor** of polynomials $p(x), q(x) \in F[x]$ if $d(x)$ evenly divides both $p(x)$ and $q(x)$; and, if for any other polynomial $d'(x)$ dividing both $p(x)$ and $q(x)$, $d'(x) \mid d(x)$. We write $d(x) = \gcd(p(x), q(x))$. Two polynomials $p(x)$ and $q(x)$ are **relatively prime** if $\gcd(p(x), q(x)) = 1$.

Proposition 17.10. *Let F be a field and suppose that $d(x)$ is a greatest common divisor of two polynomials $p(x)$ and $q(x)$ in $F[x]$. Then there exist polynomials $r(x)$ and $s(x)$ such that*

$$d(x) = r(x)p(x) + s(x)q(x).$$

Furthermore, the greatest common divisor of two polynomials is unique.

PROOF.   Let $d(x)$ be the monic polynomial of smallest degree in the set

$$S = \{f(x)p(x) + g(x)q(x) : f(x), g(x) \in F[x]\}.$$

We can write $d(x) = r(x)p(x) + s(x)q(x)$ for two polynomials $r(x)$ and $s(x)$ in $F[x]$. We need to show that $d(x)$ divides both $p(x)$ and $q(x)$. We shall first show that $d(x)$ divides

$p(x)$. By the division algorithm, there exist polynomials $a(x)$ and $b(x)$ such that $p(x) = a(x)d(x) + b(x)$, where $b(x)$ is either the zero polynomial or $\deg b(x) < \deg d(x)$. Therefore,

$$
\begin{aligned}
b(x) &= p(x) - a(x)d(x) \\
&= p(x) - a(x)(r(x)p(x) + s(x)q(x)) \\
&= p(x) - a(x)r(x)p(x) - a(x)s(x)q(x) \\
&= p(x)(1 - a(x)r(x)) + q(x)(-a(x)s(x))
\end{aligned}
$$

is a linear combination of $p(x)$ and $q(x)$ and therefore must be in S. However, $b(x)$ must be the zero polynomial since $d(x)$ was chosen to be of smallest degree; consequently, $d(x)$ divides $p(x)$. A symmetric argument shows that $d(x)$ must also divide $q(x)$; hence, $d(x)$ is a common divisor of $p(x)$ and $q(x)$.

To show that $d(x)$ is a greatest common divisor of $p(x)$ and $q(x)$, suppose that $d'(x)$ is another common divisor of $p(x)$ and $q(x)$. We will show that $d'(x) \mid d(x)$. Since $d'(x)$ is a common divisor of $p(x)$ and $q(x)$, there exist polynomials $u(x)$ and $v(x)$ such that $p(x) = u(x)d'(x)$ and $q(x) = v(x)d'(x)$. Therefore,

$$
\begin{aligned}
d(x) &= r(x)p(x) + s(x)q(x) \\
&= r(x)u(x)d'(x) + s(x)v(x)d'(x) \\
&= d'(x)[r(x)u(x) + s(x)v(x)].
\end{aligned}
$$

Since $d'(x) \mid d(x)$, $d(x)$ is a greatest common divisor of $p(x)$ and $q(x)$.

Finally, we must show that the greatest common divisor of $p(x)$ and $q(x)$ is unique. Suppose that $d'(x)$ is another greatest common divisor of $p(x)$ and $q(x)$. We have just shown that there exist polynomials $u(x)$ and $v(x)$ in $F[x]$ such that $d(x) = d'(x)[r(x)u(x) + s(x)v(x)]$. Since

$$
\deg d(x) = \deg d'(x) + \deg[r(x)u(x) + s(x)v(x)]
$$

and $d(x)$ and $d'(x)$ are both greatest common divisors, $\deg d(x) = \deg d'(x)$. Since $d(x)$ and $d'(x)$ are both monic polynomials of the same degree, it must be the case that $d(x) = d'(x)$. $\qquad \square$

Notice the similarity between the proof of Proposition 17.10 and the proof of Theorem 2.10.

17.3 Irreducible Polynomials

A nonconstant polynomial $f(x) \in F[x]$ is **irreducible** over a field F if $f(x)$ cannot be expressed as a product of two polynomials $g(x)$ and $h(x)$ in $F[x]$, where the degrees of $g(x)$ and $h(x)$ are both smaller than the degree of $f(x)$. Irreducible polynomials function as the "prime numbers" of polynomial rings.

Example 17.11. The polynomial $x^2 - 2 \in \mathbb{Q}[x]$ is irreducible since it cannot be factored any further over the rational numbers. Similarly, $x^2 + 1$ is irreducible over the real numbers.

Example 17.12. The polynomial $p(x) = x^3 + x^2 + 2$ is irreducible over $\mathbb{Z}_3[x]$. Suppose that this polynomial was reducible over $\mathbb{Z}_3[x]$. By the division algorithm there would have to be a factor of the form $x - a$, where a is some element in $\mathbb{Z}_3[x]$. Hence, it would have to be true that $p(a) = 0$. However,

$$
\begin{aligned}
p(0) &= 2 \\
p(1) &= 1
\end{aligned}
$$

$$p(2) = 2.$$

Therefore, $p(x)$ has no zeros in $\mathbb{Z}_3$ and must be irreducible.

Lemma 17.13. *Let $p(x) \in \mathbb{Q}[x]$. Then*

$$p(x) = \frac{r}{s}(a_0 + a_1 x + \cdots + a_n x^n),$$

where $r, s, a_0, \ldots, a_n$ are integers, the a_i's are relatively prime, and r and s are relatively prime.

PROOF.   Suppose that

$$p(x) = \frac{b_0}{c_0} + \frac{b_1}{c_1}x + \cdots + \frac{b_n}{c_n}x^n,$$

where the b_i's and the c_i's are integers. We can rewrite $p(x)$ as

$$p(x) = \frac{1}{c_0 \cdots c_n}(d_0 + d_1 x + \cdots + d_n x^n),$$

where $d_0, \ldots, d_n$ are integers. Let d be the greatest common divisor of $d_0, \ldots, d_n$. Then

$$p(x) = \frac{d}{c_0 \cdots c_n}(a_0 + a_1 x + \cdots + a_n x^n),$$

where $d_i = d a_i$ and the a_i's are relatively prime. Reducing $d/(c_0 \cdots c_n)$ to its lowest terms, we can write

$$p(x) = \frac{r}{s}(a_0 + a_1 x + \cdots + a_n x^n),$$

where $\gcd(r, s) = 1$.   □

Theorem 17.14 Gauss's Lemma. *Let $p(x) \in \mathbb{Z}[x]$ be a monic polynomial such that $p(x)$ factors into a product of two polynomials $\alpha(x)$ and $\beta(x)$ in $\mathbb{Q}[x]$, where the degrees of both $\alpha(x)$ and $\beta(x)$ are less than the degree of $p(x)$. Then $p(x) = a(x)b(x)$, where $a(x)$ and $b(x)$ are monic polynomials in $\mathbb{Z}[x]$ with $\deg \alpha(x) = \deg a(x)$ and $\deg \beta(x) = \deg b(x)$.*

PROOF.   By Lemma 17.13, we can assume that

$$\alpha(x) = \frac{c_1}{d_1}(a_0 + a_1 x + \cdots + a_m x^m) = \frac{c_1}{d_1}\alpha_1(x)$$
$$\beta(x) = \frac{c_2}{d_2}(b_0 + b_1 x + \cdots + b_n x^n) = \frac{c_2}{d_2}\beta_1(x),$$

where the a_i's are relatively prime and the b_i's are relatively prime. Consequently,

$$p(x) = \alpha(x)\beta(x) = \frac{c_1 c_2}{d_1 d_2}\alpha_1(x)\beta_1(x) = \frac{c}{d}\alpha_1(x)\beta_1(x),$$

where c/d is the product of c_1/d_1 and c_2/d_2 expressed in lowest terms. Hence, $dp(x) = c\alpha_1(x)\beta_1(x)$.

If $d = 1$, then $c a_m b_n = 1$ since $p(x)$ is a monic polynomial. Hence, either $c = 1$ or $c = -1$. If $c = 1$, then either $a_m = b_n = 1$ or $a_m = b_n = -1$. In the first case $p(x) = \alpha_1(x)\beta_1(x)$, where $\alpha_1(x)$ and $\beta_1(x)$ are monic polynomials with $\deg \alpha(x) = \deg \alpha_1(x)$ and $\deg \beta(x) = \deg \beta_1(x)$. In the second case $a(x) = -\alpha_1(x)$ and $b(x) = -\beta_1(x)$ are the correct monic polynomials since $p(x) = (-\alpha_1(x))(-\beta_1(x)) = a(x)b(x)$. The case in which $c = -1$ can be handled similarly.

Now suppose that $d \neq 1$. Since $\gcd(c, d) = 1$, there exists a prime p such that $p \mid d$ and $p \nmid c$. Also, since the coefficients of $\alpha_1(x)$ are relatively prime, there exists a coefficient a_i

such that $p \nmid a_i$. Similarly, there exists a coefficient b_j of $\beta_1(x)$ such that $p \nmid b_j$. Let $\alpha_1'(x)$ and $\beta_1'(x)$ be the polynomials in $\mathbb{Z}_p[x]$ obtained by reducing the coefficients of $\alpha_1(x)$ and $\beta_1(x)$ modulo p. Since $p \mid d$, $\alpha_1'(x)\beta_1'(x) = 0$ in $\mathbb{Z}_p[x]$. However, this is impossible since neither $\alpha_1'(x)$ nor $\beta_1'(x)$ is the zero polynomial and $\mathbb{Z}_p[x]$ is an integral domain. Therefore, $d = 1$ and the theorem is proven. $\qquad\square$

Corollary 17.15. *Let* $p(x) = x^n + a_{n-1}x^{n-1} + \cdots + a_0$ *be a polynomial with coefficients in* $\mathbb{Z}$ *and* $a_0 \neq 0$. *If* $p(x)$ *has a zero in* $\mathbb{Q}$, *then* $p(x)$ *also has a zero* α *in* $\mathbb{Z}$. *Furthermore,* α *divides* a_0.

PROOF. Let $p(x)$ have a zero $a \in \mathbb{Q}$. Then $p(x)$ must have a linear factor $x - a$. By Gauss's Lemma, $p(x)$ has a factorization with a linear factor in $\mathbb{Z}[x]$. Hence, for some $\alpha \in \mathbb{Z}$

$$p(x) = (x - \alpha)(x^{n-1} + \cdots - a_0/\alpha).$$

Thus $a_0/\alpha \in \mathbb{Z}$ and so $\alpha \mid a_0$. $\qquad\square$

Example 17.16. Let $p(x) = x^4 - 2x^3 + x + 1$. We shall show that $p(x)$ is irreducible over $\mathbb{Q}[x]$. Assume that $p(x)$ is reducible. Then either $p(x)$ has a linear factor, say $p(x) = (x - \alpha)q(x)$, where $q(x)$ is a polynomial of degree three, or $p(x)$ has two quadratic factors.

If $p(x)$ has a linear factor in $\mathbb{Q}[x]$, then it has a zero in $\mathbb{Z}$. By Corollary 17.15, any zero must divide 1 and therefore must be ± 1; however, $p(1) = 1$ and $p(-1) = 3$. Consequently, we have eliminated the possibility that $p(x)$ has any linear factors.

Therefore, if $p(x)$ is reducible it must factor into two quadratic polynomials, say

$$\begin{aligned} p(x) &= (x^2 + ax + b)(x^2 + cx + d) \\ &= x^4 + (a + c)x^3 + (ac + b + d)x^2 + (ad + bc)x + bd, \end{aligned}$$

where each factor is in $\mathbb{Z}[x]$ by Gauss's Lemma. Hence,

$$\begin{aligned} a + c &= -2 \\ ac + b + d &= 0 \\ ad + bc &= 1 \\ bd &= 1. \end{aligned}$$

Since $bd = 1$, either $b = d = 1$ or $b = d = -1$. In either case $b = d$ and so

$$ad + bc = b(a + c) = 1.$$

Since $a + c = -2$, we know that $-2b = 1$. This is impossible since b is an integer. Therefore, $p(x)$ must be irreducible over $\mathbb{Q}$.

Theorem 17.17 Eisenstein's Criterion. *Let* p *be a prime and suppose that*

$$f(x) = a_n x^n + \cdots + a_0 \in \mathbb{Z}[x].$$

If $p \mid a_i$ *for* $i = 0, 1, \ldots, n - 1$, *but* $p \nmid a_n$ *and* $p^2 \nmid a_0$, *then* $f(x)$ *is irreducible over* $\mathbb{Q}$.

PROOF. By Gauss's Lemma, we need only show that $f(x)$ does not factor into polynomials of lower degree in $\mathbb{Z}[x]$. Let

$$f(x) = (b_r x^r + \cdots + b_0)(c_s x^s + \cdots + c_0)$$

be a factorization in $\mathbb{Z}[x]$, with b_r and c_s not equal to zero and $r, s < n$. Since p^2 does not divide $a_0 = b_0 c_0$, either b_0 or c_0 is not divisible by p. Suppose that $p \nmid b_0$ and $p \mid c_0$. Since

$p \nmid a_n$ and $a_n = b_r c_s$, neither b_r nor c_s is divisible by p. Let m be the smallest value of k such that $p \nmid c_k$. Then

$$a_m = b_0 c_m + b_1 c_{m-1} + \cdots + b_m c_0$$

is not divisible by p, since each term on the right-hand side of the equation is divisible by p except for $b_0 c_m$. Therefore, $m = n$ since a_i is divisible by p for $m < n$. Hence, $f(x)$ cannot be factored into polynomials of lower degree and therefore must be irreducible. $\square$

Example 17.18. The polynomial

$$f(x) = 16x^5 - 9x^4 + 3x^2 + 6x - 21$$

is easily seen to be irreducible over $\mathbb{Q}$ by Eisenstein's Criterion if we let $p = 3$.

Eisenstein's Criterion is more useful in constructing irreducible polynomials of a certain degree over $\mathbb{Q}$ than in determining the irreducibility of an arbitrary polynomial in $\mathbb{Q}[x]$: given an arbitrary polynomial, it is not very likely that we can apply Eisenstein's Criterion. The real value of Theorem 17.17 is that we now have an easy method of generating irreducible polynomials of any degree.

Ideals in $F[x]$

Let F be a field. Recall that a principal ideal in $F[x]$ is an ideal $\langle p(x) \rangle$ generated by some polynomial $p(x)$; that is,

$$\langle p(x) \rangle = \{ p(x)q(x) : q(x) \in F[x] \}.$$

Example 17.19. The polynomial x^2 in $F[x]$ generates the ideal $\langle x^2 \rangle$ consisting of all polynomials with no constant term or term of degree 1.

Theorem 17.20. *If F is a field, then every ideal in $F[x]$ is a principal ideal.*

PROOF.   Let I be an ideal of $F[x]$. If I is the zero ideal, the theorem is easily true. Suppose that I is a nontrivial ideal in $F[x]$, and let $p(x) \in I$ be a nonzero element of minimal degree. If $\deg p(x) = 0$, then $p(x)$ is a nonzero constant and 1 must be in I. Since 1 generates all of $F[x]$, $\langle 1 \rangle = I = F[x]$ and I is again a principal ideal.

Now assume that $\deg p(x) \geq 1$ and let $f(x)$ be any element in I. By the division algorithm there exist $q(x)$ and $r(x)$ in $F[x]$ such that $f(x) = p(x)q(x) + r(x)$ and $\deg r(x) < \deg p(x)$. Since $f(x), p(x) \in I$ and I is an ideal, $r(x) = f(x) - p(x)q(x)$ is also in I. However, since we chose $p(x)$ to be of minimal degree, $r(x)$ must be the zero polynomial. Since we can write any element $f(x)$ in I as $p(x)q(x)$ for some $q(x) \in F[x]$, it must be the case that $I = \langle p(x) \rangle$. $\square$

Example 17.21. It is not the case that every ideal in the ring $F[x, y]$ is a principal ideal. Consider the ideal of $F[x, y]$ generated by the polynomials x and y. This is the ideal of $F[x, y]$ consisting of all polynomials with no constant term. Since both x and y are in the ideal, no single polynomial can generate the entire ideal.

Theorem 17.22. *Let F be a field and suppose that $p(x) \in F[x]$. Then the ideal generated by $p(x)$ is maximal if and only if $p(x)$ is irreducible.*

PROOF.   Suppose that $p(x)$ generates a maximal ideal of $F[x]$. Then $\langle p(x) \rangle$ is also a prime ideal of $F[x]$. Since a maximal ideal must be properly contained inside $F[x]$, $p(x)$ cannot be a constant polynomial. Let us assume that $p(x)$ factors into two polynomials of lesser degree, say $p(x) = f(x)g(x)$. Since $\langle p(x) \rangle$ is a prime ideal one of these factors, say $f(x)$, is

in $\langle p(x) \rangle$ and therefore be a multiple of $p(x)$. But this would imply that $\langle p(x) \rangle \subset \langle f(x) \rangle$, which is impossible since $\langle p(x) \rangle$ is maximal.

Conversely, suppose that $p(x)$ is irreducible over $F[x]$. Let I be an ideal in $F[x]$ containing $\langle p(x) \rangle$. By Theorem 17.20, I is a principal ideal; hence, $I = \langle f(x) \rangle$ for some $f(x) \in F[x]$. Since $p(x) \in I$, it must be the case that $p(x) = f(x)g(x)$ for some $g(x) \in F[x]$. However, $p(x)$ is irreducible; hence, either $f(x)$ or $g(x)$ is a constant polynomial. If $f(x)$ is constant, then $I = F[x]$ and we are done. If $g(x)$ is constant, then $f(x)$ is a constant multiple of I and $I = \langle p(x) \rangle$. Thus, there are no proper ideals of $F[x]$ that properly contain $\langle p(x) \rangle$. $\qquad\square$

Historical Note

Throughout history, the solution of polynomial equations has been a challenging problem. The Babylonians knew how to solve the equation $ax^2 + bx + c = 0$. Omar Khayyam (1048–1131) devised methods of solving cubic equations through the use of geometric constructions and conic sections. The algebraic solution of the general cubic equation $ax^3 + bx^2 + cx + d = 0$ was not discovered until the sixteenth century. An Italian mathematician, Luca Pacioli (ca. 1445–1509), wrote in *Summa de Arithmetica* that the solution of the cubic was impossible. This was taken as a challenge by the rest of the mathematical community.

Scipione del Ferro (1465–1526), of the University of Bologna, solved the "depressed cubic,"

$$ax^3 + cx + d = 0.$$

He kept his solution an absolute secret. This may seem surprising today, when mathematicians are usually very eager to publish their results, but in the days of the Italian Renaissance secrecy was customary. Academic appointments were not easy to secure and depended on the ability to prevail in public contests. Such challenges could be issued at any time. Consequently, any major new discovery was a valuable weapon in such a contest. If an opponent presented a list of problems to be solved, del Ferro could in turn present a list of depressed cubics. He kept the secret of his discovery throughout his life, passing it on only on his deathbed to his student Antonio Fior (ca. 1506–?).

Although Fior was not the equal of his teacher, he immediately issued a challenge to Niccolo Fontana (1499–1557). Fontana was known as Tartaglia (the Stammerer). As a youth he had suffered a blow from the sword of a French soldier during an attack on his village. He survived the savage wound, but his speech was permanently impaired. Tartaglia sent Fior a list of 30 various mathematical problems; Fior countered by sending Tartaglia a list of 30 depressed cubics. Tartaglia would either solve all 30 of the problems or absolutely fail. After much effort Tartaglia finally succeeded in solving the depressed cubic and defeated Fior, who faded into obscurity.

At this point another mathematician, Gerolamo Cardano (1501–1576), entered the story. Cardano wrote to Tartaglia, begging him for the solution to the depressed cubic. Tartaglia refused several of his requests, then finally revealed the solution to Cardano after the latter swore an oath not to publish the secret or to pass it on to anyone else. Using the knowledge that he had obtained from Tartaglia, Cardano eventually solved the general cubic

$$ax^3 + bx^2 + cx + d = 0.$$

Cardano shared the secret with his student, Ludovico Ferrari (1522–1565), who solved the general quartic equation,

$$ax^4 + bx^3 + cx^2 + dx + e = 0.$$

In 1543, Cardano and Ferrari examined del Ferro's papers and discovered that he had also solved the depressed cubic. Cardano felt that this relieved him of his obligation to Tartaglia, so he proceeded to publish the solutions in *Ars Magna* (1545), in which he gave credit to del Ferro for solving the special case of the cubic. This resulted in a bitter dispute between Cardano and Tartaglia, who published the story of the oath a year later.

17.4   Exercises

1.   List all of the polynomials of degree 3 or less in $\mathbb{Z}_2[x]$.

2.   Compute each of the following.

(a) $(5x^2 + 3x - 4) + (4x^2 - x + 9)$ in $\mathbb{Z}_{12}$

(b) $(5x^2 + 3x - 4)(4x^2 - x + 9)$ in $\mathbb{Z}_{12}$

(c) $(7x^3 + 3x^2 - x) + (6x^2 - 8x + 4)$ in $\mathbb{Z}_9$

(d) $(3x^2 + 2x - 4) + (4x^2 + 2)$ in $\mathbb{Z}_5$

(e) $(3x^2 + 2x - 4)(4x^2 + 2)$ in $\mathbb{Z}_5$

(f) $(5x^2 + 3x - 2)^2$ in $\mathbb{Z}_{12}$

3.   Use the division algorithm to find $q(x)$ and $r(x)$ such that $a(x) = q(x)b(x) + r(x)$ with $\deg r(x) < \deg b(x)$ for each of the following pairs of polynomials.

(a) $a(x) = 5x^3 + 6x^2 - 3x + 4$ and $b(x) = x - 2$ in $\mathbb{Z}_7[x]$

(b) $a(x) = 6x^4 - 2x^3 + x^2 - 3x + 1$ and $b(x) = x^2 + x - 2$ in $\mathbb{Z}_7[x]$

(c) $a(x) = 4x^5 - x^3 + x^2 + 4$ and $b(x) = x^3 - 2$ in $\mathbb{Z}_5[x]$

(d) $a(x) = x^5 + x^3 - x^2 - x$ and $b(x) = x^3 + x$ in $\mathbb{Z}_2[x]$

4.   Find the greatest common divisor of each of the following pairs $p(x)$ and $q(x)$ of polynomials. If $d(x) = \gcd(p(x), q(x))$, find two polynomials $a(x)$ and $b(x)$ such that $a(x)p(x) + b(x)q(x) = d(x)$.

(a) $p(x) = x^3 - 6x^2 + 14x - 15$ and $q(x) = x^3 - 8x^2 + 21x - 18$, where $p(x), q(x) \in \mathbb{Q}[x]$

(b) $p(x) = x^3 + x^2 - x + 1$ and $q(x) = x^3 + x - 1$, where $p(x), q(x) \in \mathbb{Z}_2[x]$

(c) $p(x) = x^3 + x^2 - 4x + 4$ and $q(x) = x^3 + 3x - 2$, where $p(x), q(x) \in \mathbb{Z}_5[x]$

(d) $p(x) = x^3 - 2x + 4$ and $q(x) = 4x^3 + x + 3$, where $p(x), q(x) \in \mathbb{Q}[x]$

5.   Find all of the zeros for each of the following polynomials.

(a) $5x^3 + 4x^2 - x + 9$ in $\mathbb{Z}_{12}$ (c) $5x^4 + 2x^2 - 3$ in $\mathbb{Z}_7$

(b) $3x^3 - 4x^2 - x + 4$ in $\mathbb{Z}_5$ (d) $x^3 + x + 1$ in $\mathbb{Z}_2$

6.   Find all of the units in $\mathbb{Z}[x]$.

7.   Find a unit $p(x)$ in $\mathbb{Z}_4[x]$ such that $\deg p(x) > 1$.

8.   Which of the following polynomials are irreducible over $\mathbb{Q}[x]$?

(a) $x^4 - 2x^3 + 2x^2 + x + 4$

(c) $3x^5 - 4x^3 - 6x^2 + 6$

(b) $x^4 - 5x^3 + 3x - 2$

(d) $5x^5 - 6x^4 - 3x^2 + 9x - 15$

9.   Find all of the irreducible polynomials of degrees 2 and 3 in $\mathbb{Z}_2[x]$.

10.   Give two different factorizations of $x^2 + x + 8$ in $\mathbb{Z}_{10}[x]$.

11.   Prove or disprove: There exists a polynomial $p(x)$ in $\mathbb{Z}_6[x]$ of degree n with more than n distinct zeros.

12.   If F is a field, show that $F[x_1, \ldots, x_n]$ is an integral domain.

13.   Show that the division algorithm does not hold for $\mathbb{Z}[x]$. Why does it fail?

14.   Prove or disprove: $x^p + a$ is irreducible for any $a \in \mathbb{Z}_p$, where p is prime.

15.   Let $f(x)$ be irreducible in $F[x]$, where F is a field. If $f(x) \mid p(x)q(x)$, prove that either $f(x) \mid p(x)$ or $f(x) \mid q(x)$.

16.   Suppose that R and S are isomorphic rings. Prove that $R[x] \cong S[x]$.

17.   Let F be a field and $a \in F$. If $p(x) \in F[x]$, show that $p(a)$ is the remainder obtained when $p(x)$ is divided by $x - a$.

18. **The Rational Root Theorem.**   Let

$$p(x) = a_n x^n + a_{n-1} x^{n-1} + \cdots + a_0 \in \mathbb{Z}[x],$$

where $a_n \neq 0$. Prove that if $p(r/s) = 0$, where $\gcd(r, s) = 1$, then $r \mid a_0$ and $s \mid a_n$.

19.   Let $\mathbb{Q}^*$ be the multiplicative group of positive rational numbers. Prove that $\mathbb{Q}^*$ is isomorphic to $(\mathbb{Z}[x], +)$.

20. **Cyclotomic Polynomials.**   The polynomial

$$\Phi_n(x) = \frac{x^n - 1}{x - 1} = x^{n-1} + x^{n-2} + \cdots + x + 1$$

is called the ***cyclotomic polynomial.*** Show that $\Phi_p(x)$ is irreducible over $\mathbb{Q}$ for any prime p.

21.   If F is a field, show that there are infinitely many irreducible polynomials in $F[x]$.

22.   Let R be a commutative ring with identity. Prove that multiplication is commutative in $R[x]$.

23.   Let R be a commutative ring with identity. Prove that multiplication is distributive in $R[x]$.

24.   Show that $x^p - x$ has p distinct zeros in $\mathbb{Z}_p$, for any prime p. Conclude that

$$x^p - x = x(x - 1)(x - 2) \cdots (x - (p - 1)).$$

25.   Let F be a field and $f(x) = a_0 + a_1 x + \cdots + a_n x^n$ be in $F[x]$. Define $f'(x) = a_1 + 2a_2 x + \cdots + n a_n x^{n-1}$ to be the ***derivative*** of $f(x)$.

(a) Prove that

$$(f + g)'(x) = f'(x) + g'(x).$$

Conclude that we can define a homomorphism of abelian groups $D : F[x] \to F[x]$ by $D(f(x)) = f'(x)$.

(b) Calculate the kernel of D if char $F = 0$.

(c) Calculate the kernel of D if char $F = p$.

(d) Prove that
$$(fg)'(x) = f'(x)g(x) + f(x)g'(x).$$

(e) Suppose that we can factor a polynomial $f(x) \in F[x]$ into linear factors, say
$$f(x) = a(x - a_1)(x - a_2) \cdots (x - a_n).$$

Prove that $f(x)$ has no repeated factors if and only if $f(x)$ and $f'(x)$ are relatively prime.

26. Let F be a field. Show that $F[x]$ is never a field.

27. Let R be an integral domain. Prove that $R[x_1, \ldots, x_n]$ is an integral domain.

28. Let R be a commutative ring with identity. Show that $R[x]$ has a subring R' isomorphic to R.

29. Let $p(x)$ and $q(x)$ be polynomials in $R[x]$, where R is a commutative ring with identity. Prove that $\deg(p(x) + q(x)) \leq \max(\deg p(x), \deg q(x))$.

17.5 Additional Exercises: Solving the Cubic and Quartic Equations

1. Solve the general quadratic equation
$$ax^2 + bx + c = 0$$

to obtain
$$x = \frac{-b \pm \sqrt{b^2 - 4ac}}{2a}.$$

The **discriminant** of the quadratic equation $\Delta = b^2 - 4ac$ determines the nature of the solutions of the equation. If $\Delta > 0$, the equation has two distinct real solutions. If $\Delta = 0$, the equation has a single repeated real root. If $\Delta < 0$, there are two distinct imaginary solutions.

2. Show that any cubic equation of the form
$$x^3 + bx^2 + cx + d = 0$$

can be reduced to the form $y^3 + py + q = 0$ by making the substitution $x = y - b/3$.

3. Prove that the cube roots of 1 are given by

$$\omega = \frac{-1 + i\sqrt{3}}{2}$$
$$\omega^2 = \frac{-1 - i\sqrt{3}}{2}$$
$$\omega^3 = 1.$$

4. Make the substitution
$$y = z - \frac{p}{3z}$$

for y in the equation $y^3 + py + q = 0$ and obtain two solutions A and B for z^3.

5. Show that the product of the solutions obtained in (4) is $-p^3/27$, deducing that $\sqrt[3]{AB} = -p/3$.

6. Prove that the possible solutions for z in (4) are given by

$$\sqrt[3]{A}, \quad \omega\sqrt[3]{A}, \quad \omega^2\sqrt[3]{A}, \quad \sqrt[3]{B}, \quad \omega\sqrt[3]{B}, \quad \omega^2\sqrt[3]{B}$$

and use this result to show that the three possible solutions for y are

$$\omega^i \sqrt[3]{-\frac{q}{2} + \sqrt{\frac{p^3}{27} + \frac{q^2}{4}}} + \omega^{2i} \sqrt[3]{-\frac{q}{2} - \sqrt{\frac{p^3}{27} + \frac{q^2}{4}}},$$

where $i = 0, 1, 2$.

7. The **discriminant** of the cubic equation is

$$\Delta = \frac{p^3}{27} + \frac{q^2}{4}.$$

Show that $y^3 + py + q = 0$

(a) has three real roots, at least two of which are equal, if $\Delta = 0$.

(b) has one real root and two conjugate imaginary roots if $\Delta > 0$.

(c) has three distinct real roots if $\Delta < 0$.

8. Solve the following cubic equations.

(a) $x^3 - 4x^2 + 11x + 30 = 0$

(b) $x^3 - 3x + 5 = 0$

(c) $x^3 - 3x + 2 = 0$

(d) $x^3 + x + 3 = 0$

9. Show that the general quartic equation

$$x^4 + ax^3 + bx^2 + cx + d = 0$$

can be reduced to

$$y^4 + py^2 + qy + r = 0$$

by using the substitution $x = y - a/4$.

10. Show that

$$\left(y^2 + \frac{1}{2}z\right)^2 = (z - p)y^2 - qy + \left(\frac{1}{4}z^2 - r\right).$$

11. Show that the right-hand side of Exercise 17.5.10 can be put in the form $(my + k)^2$ if and only if

$$q^2 - 4(z - p)\left(\frac{1}{4}z^2 - r\right) = 0.$$

12. From Exercise 17.5.11 obtain the **resolvent cubic equation**

$$z^3 - pz^2 - 4rz + (4pr - q^2) = 0.$$

Solving the resolvent cubic equation, put the equation found in Exercise 17.5.10 in the form

$$\left(y^2 + \frac{1}{2}z\right)^2 = (my + k)^2$$

to obtain the solution of the quartic equation.

13. Use this method to solve the following quartic equations.

(a) $x^4 - x^2 - 3x + 2 = 0$

(b) $x^4 + x^3 - 7x^2 - x + 6 = 0$

(c) $x^4 - 2x^2 + 4x - 3 = 0$

(d) $x^4 - 4x^3 + 3x^2 - 5x + 2 = 0$

17.6   Sage

Sage is particularly adept at building, analyzing and manipulating polynomial rings. We have seen some of this in the previous chapter. Let's begin by creating three polynomial rings and checking some of their basic properties. There are several ways to construct polynomial rings, but the syntax used here is the most straightforward.

Polynomial Rings and their Elements

```
R.<x> = Integers(8)[]; R
```

Univariate Polynomial Ring **in** x over Ring of integers modulo 8

```
S.<y> = ZZ[]; S
```

Univariate Polynomial Ring **in** y over Integer Ring

```
T.<z> = QQ[]; T
```

Univariate Polynomial Ring **in** z over Rational Field

Basic properties of rings are availble for these examples.

```
R.is_finite()
```

False

```
R.is_integral_domain()
```

False

```
S.is_integral_domain()
```

True

```
T.is_field()
```

False

```
R.characteristic()
```

8

```
T.characteristic()
```

0

With the construction syntax used above, the variables can be used to create elements of the polynomial ring without explicit coercion (though we need to be careful about constant polynomials).

```
y in S
```

True

```
x in S
```

False

```
q = (3/2) + (5/4)*z^2
q in T
```

True

```
3 in S
```

True

```
r = 3
r.parent()
```

Integer Ring

```
s = 3*y^0
s.parent()
```

Univariate Polynomial Ring in y over Integer Ring

Polynomials can be evaluated like they are functions, so we can mimic the evaluation homomorphism.

```
p = 3 + 5*x + 2*x^2
p.parent()
```

Univariate Polynomial Ring in x over Ring of integers modulo 8

```
p(1)
```

2

```
[p(t) for t in Integers(8)]
```

[3, 2, 5, 4, 7, 6, 1, 0]

Notice that p is a degree two polynomial, yet through a brute-force examination we see that the polynomial only has one root, contrary to our usual expectations. It can be even more unusual.

```
q = 4*x^2+4*x
[q(t) for t in Integers(8)]
```

[0, 0, 0, 0, 0, 0, 0, 0]

Sage can create and manipulate rings of polynomials in more than one variable, though we will not have much occasion to use this functionality in this course.

```
M.<s, t> = QQ[]; M
```

Multivariate Polynomial Ring in s, t over Rational Field

Irreducible Polynomials

Whether or not a polynomial factors, taking into consideration the ring used for its coefficients, is an important topic in this chapter and many of the following chapters. Sage can factor, and determine irreducibility, over the integers, the rationals, and finite fields.

First, over the rationals.

```
R.<x> = QQ[]
p = 1/4*x^4 - x^3 + x^2 - x - 1/2
p.is_irreducible()
```

True

```
p.factor()
```

(1/4) * (x^4 - 4*x^3 + 4*x^2 - 4*x - 2)

```
q = 2*x^5 + 5/2*x^4 + 3/4*x^3 - 25/24*x^2 - x - 1/2
q.is_irreducible()
```

False

```
q.factor()
```

(2) * (x^2 + 3/2*x + 3/4) * (x^3 - 1/4*x^2 - 1/3)

Factoring over the integers is really no different than factoring over the rationals. This is the content of Theorem 17.14 — finding a factorization over the integers can be converted to finding a factorization over the rationals. So it is with Sage, there is little difference between working over the rationals and the integers. It is a little different working over a finite field. Commentary follows.

```
F.<a> = FiniteField(5^2)
S.<y> = F[]
p = 2*y^5 + 2*y^4 + 4*y^3 + 2*y^2 + 3*y + 1
p.is_irreducible()
```

True

```
p.factor()
```

(2) * (y^5 + y^4 + 2*y^3 + y^2 + 4*y + 3)

```
q = 3*y^4+2*y^3-y+4; q.factor()
```

(3) * (y^2 + (a + 4)*y + 2*a + 3) * (y^2 + 4*a*y + 3*a)

```
r = y^4+2*y^3+3*y^2+4; r.factor()
```

(y + 4) * (y^3 + 3*y^2 + y + 1)

```
s = 3*y^4+2*y^3-y+3; s.factor()
```

(3) * (y + 1) * (y + 3) * (y + 2*a + 4) * (y + 3*a + 1)

To check these factorizations, we need to compute in the finite field, F, and so we need to know how the symbol a behaves. This symbol is considered as a root of a degree two polynomial over the integers mod 5, which we can get with the .modulus() method.

```
F.modulus()
```

```
x^2 + 4*x + 2
```

So $a^2 + 4a + 2 = 0$, or $a^2 = -4a - 3 = a + 2$. So when checking the factorizations, anytime you see an a^2 you can replace it by $a + 2$. Notice that by Corollary 17.8 we could find the one linear factor of r, and the four linear factors of s, through a brute-force search for roots. This is feasible because the field is finite.

```
[t for t in F if r(t)==0]
```

```
[1]
```

```
[t for t in F if s(t)==0]
```

```
[2, 3*a + 1, 4, 2*a + 4]
```

However, q factors into a pair of degree 2 polynomials, so no amount of testing for roots will discover a factor.

With Eisenstein's Criterion, we can create irreducible polynomials, such as in Example 17.18.

```
W.<w> = QQ[]
p = 16*w^5 - 9*w^4 +3*w^2 + 6*w -21
p.is_irreducible()
```

```
True
```

Over the field $\mathbb{Z}_p$, the field of integers mod a prime p, Conway polynomials are canonical choices of a polynomial of degree n that is irreducible over $\mathbb{Z}_p$. See the exercises for more about these polynomials.

Polynomials over Fields

If F is a field, then every ideal of $F[x]$ is principal (Theorem 17.20). Nothing stops you from giving Sage two (or more) generators to construct an ideal, but Sage will determine the element to use in a description of the ideal as a principal ideal.

```
W.<w> = QQ[]
r = -w^5 + 5*w^4 - 4*w^3 + 14*w^2 - 67*w + 17
s = 3*w^5 - 14*w^4 + 12*w^3 - 6*w^2 + w
S = W.ideal(r, s)
S
```

```
Principal ideal (w^2 - 4*w + 1) of
Univariate Polynomial Ring in w over Rational Field
```

```
(w^2)*r + (3*w-6)*s in S
```

```
True
```

Theorem 17.22 is the key fact that allows us to easily construct finite fields. Here is a construction of a finite field of order $7^5 = 16\,807$. All we need is a polynomial of degree 5 that is irreducible over $\mathbb{Z}_7$.

```
F = Integers(7)
R.<x> = F[]
p = x^5+ x + 4
p.is_irreducible()
```

True

```
id = R.ideal(p)
Q = R.quotient(id); Q
```

Univariate Quotient Polynomial Ring **in** xbar over
Ring of integers modulo 7 with modulus x^5 + x + 4

```
Q.is_field()
```

True

```
Q.order() == 7^5
```

True

The symbol xbar is a generator of the field, but right now it is not accessible. xbar is the
coset $x + \langle x^5 + x + 4 \rangle$. A better construction would include specifying this generator.

```
Q.gen(0)
```

xbar

```
Q.<t> = R.quotient(id); Q
```

Univariate Quotient Polynomial Ring **in** t over
Ring of integers modulo 7 with modulus x^5 + x + 4

```
t^5 + t + 4
```

0

```
t^5 == -(t+4)
```

True

```
t^5
```

6*t + 3

```
(3*t^3 + t + 5)*(t^2 + 4*t + 2)
```

5*t^4 + 2*t^2 + 5*t + 5

```
a = 3*t^4 - 6*t^3 + 3*t^2 + 5*t + 2
ainv = a^-1; ainv
```

6*t^4 + 5*t^2 + 4

```
a*ainv
```

1

17.7 Sage Exercises

1. Consider the polynomial $x^3 - 3x + 4$. Compute the most thorough factorization of this polynomial over each of the following fields: (a) the finite field $\mathbb{Z}_5$, (b) a finite field with 125 elements, (c) the rationals, (d) the real numbers and (e) the complex numbers. To do this, build the appropriate polynomial ring, and construct the polynomial as a member of this ring, and use the `.factor()` method.

2. "Conway polynomials" are irreducible polynomials over $\mathbb{Z}_p$ that Sage (and other software) uses to build maximal ideals in polynomial rings, and thus quotient rings that are fields. Roughly speaking, they are "canonical" choices for each degree and each prime. The command `conway_polynomial(p, n)` will return a database entry that is an irreducible polynomial of degree n over $\mathbb{Z}_p$.

Execute the command `conway_polynomial(5, 4)` to obtain an allegedly irreducible polynomial of degree 4 over $\mathbb{Z}_5$: $p = x^4 + 4x^2 + 4x + 2$. Construct the right polynomial ring (i.e., in the indeterminate x) and verify that `p` is really an element of your polynomial ring.

First determine that p has no linear factors. The only possibility left is that p factors as two quadratic polynomials over $\mathbb{Z}_5$. Use a list comprehension with *three* `for` statements to create *every* possible quadratic polynomial over $\mathbb{Z}_5$. Now use this list to create every possible product of two quadratic polynomials and check to see if `p` is in this list.

More on Conway polynomials is available at Frank Lübeck's site.

3. Construct a finite field of order 729 as a quotient of a polynomial ring by a principal ideal generated with a Conway polynomial.

4. Define the polynomials $p = x^3 + 2x^2 + 2x + 4$ and $q = x^4 + 2x^2$ as polynomials with coefficients from the integers. Compute `gcd(p, q)` and verify that the result divides both `p` and `q` (just form a fraction in Sage and see that it simplifies cleanly, or use the `.quo_rem()` method).

Proposition 17.10 says there are polynomials $r(x)$ and $s(x)$ such that the greatest common divisor equals $r(x)p(x) + s(x)q(x)$, *if the coefficients come from a field.* Since here we have two polynomials over the integers, investigate the results returned by Sage for the extended gcd, `xgcd(p, q)`. In particular, show that the first result of the returned triple is a multiple of the gcd. Then verify the "linear combination" property of the result.

5. For a polynomial ring over a field, every ideal is principal. Begin with the ring of polynomials over the rationals. Experiment with constructing ideals using two generators and then see that Sage converts the ideal to a principal ideal with a single generator. (You can get this generator with the ideal method `.gen()`.) Can you explain how this single generator is computed?

<div align="center">

18

Integral Domains

</div>

One of the most important rings we study is the ring of integers. It was our first example of an algebraic structure: the first polynomial ring that we examined was $\mathbb{Z}[x]$. We also know that the integers sit naturally inside the field of rational numbers, $\mathbb{Q}$. The ring of integers is the model for all integral domains. In this chapter we will examine integral domains in general, answering questions about the ideal structure of integral domains, polynomial rings over integral domains, and whether or not an integral domain can be embedded in a field.

18.1 Fields of Fractions

Every field is also an integral domain; however, there are many integral domains that are not fields. For example, the integers $\mathbb{Z}$ form an integral domain but not a field. A question that naturally arises is how we might associate an integral domain with a field. There is a natural way to construct the rationals $\mathbb{Q}$ from the integers: the rationals can be represented as formal quotients of two integers. The rational numbers are certainly a field. In fact, it can be shown that the rationals are the smallest field that contains the integers. Given an integral domain D, our question now becomes how to construct a smallest field F containing D. We will do this in the same way as we constructed the rationals from the integers.

An element $p/q \in \mathbb{Q}$ is the quotient of two integers p and q; however, different pairs of integers can represent the same rational number. For instance, $1/2 = 2/4 = 3/6$. We know that

$$\frac{a}{b} = \frac{c}{d}$$

if and only if $ad = bc$. A more formal way of considering this problem is to examine fractions in terms of equivalence relations. We can think of elements in $\mathbb{Q}$ as ordered pairs in $\mathbb{Z} \times \mathbb{Z}$. A quotient p/q can be written as (p, q). For instance, $(3, 7)$ would represent the fraction $3/7$. However, there are problems if we consider all possible pairs in $\mathbb{Z} \times \mathbb{Z}$. There is no fraction $5/0$ corresponding to the pair $(5, 0)$. Also, the pairs $(3, 6)$ and $(2, 4)$ both represent the fraction $1/2$. The first problem is easily solved if we require the second coordinate to be nonzero. The second problem is solved by considering two pairs (a, b) and (c, d) to be equivalent if $ad = bc$.

If we use the approach of ordered pairs instead of fractions, then we can study integral domains in general. Let D be any integral domain and let

$$S = \{(a, b) : a, b \in D \text{ and } b \neq 0\}.$$

Define a relation on S by $(a, b) \sim (c, d)$ if $ad = bc$.

Lemma 18.1. *The relation $\sim$ between elements of S is an equivalence relation.*

PROOF. Since D is commutative, $ab = ba$; hence, $\sim$ is reflexive on D. Now suppose that $(a, b) \sim (c, d)$. Then $ad = bc$ or $cb = da$. Therefore, $(c, d) \sim (a, b)$ and the relation is symmetric. Finally, to show that the relation is transitive, let $(a, b) \sim (c, d)$ and $(c, d) \sim (e, f)$. In this case $ad = bc$ and $cf = de$. Multiplying both sides of $ad = bc$ by f yields

$$afd = adf = bcf = bde = bed.$$

Since D is an integral domain, we can deduce that $af = be$ or $(a, b) \sim (e, f)$. □

We will denote the set of equivalence classes on S by F_D. We now need to define the operations of addition and multiplication on F_D. Recall how fractions are added and multiplied in $\mathbb{Q}$:

$$\frac{a}{b} + \frac{c}{d} = \frac{ad + bc}{bd};$$
$$\frac{a}{b} \cdot \frac{c}{d} = \frac{ac}{bd}.$$

It seems reasonable to define the operations of addition and multiplication on F_D in a similar manner. If we denote the equivalence class of $(a, b) \in S$ by $[a, b]$, then we are led to define the operations of addition and multiplication on F_D by

$$[a, b] + [c, d] = [ad + bc, bd]$$

and

$$[a, b] \cdot [c, d] = [ac, bd],$$

respectively. The next lemma demonstrates that these operations are independent of the choice of representatives from each equivalence class.

Lemma 18.2. *The operations of addition and multiplication on F_D are well-defined.*

PROOF. We will prove that the operation of addition is well-defined. The proof that multiplication is well-defined is left as an exercise. Let $[a_1, b_1] = [a_2, b_2]$ and $[c_1, d_1] = [c_2, d_2]$. We must show that

$$[a_1 d_1 + b_1 c_1, b_1 d_1] = [a_2 d_2 + b_2 c_2, b_2 d_2]$$

or, equivalently, that

$$(a_1 d_1 + b_1 c_1)(b_2 d_2) = (b_1 d_1)(a_2 d_2 + b_2 c_2).$$

Since $[a_1, b_1] = [a_2, b_2]$ and $[c_1, d_1] = [c_2, d_2]$, we know that $a_1 b_2 = b_1 a_2$ and $c_1 d_2 = d_1 c_2$. Therefore,

$$\begin{aligned}
(a_1 d_1 + b_1 c_1)(b_2 d_2) &= a_1 d_1 b_2 d_2 + b_1 c_1 b_2 d_2 \\
&= a_1 b_2 d_1 d_2 + b_1 b_2 c_1 d_2 \\
&= b_1 a_2 d_1 d_2 + b_1 b_2 d_1 c_2 \\
&= (b_1 d_1)(a_2 d_2 + b_2 c_2). \quad \square
\end{aligned}$$

Lemma 18.3. *The set of equivalence classes of S, F_D, under the equivalence relation $\sim$, together with the operations of addition and multiplication defined by*

$$[a, b] + [c, d] = [ad + bc, bd]$$
$$[a, b] \cdot [c, d] = [ac, bd],$$

is a field.

PROOF.   The additive and multiplicative identities are $[0, 1]$ and $[1, 1]$, respectively. To show that $[0, 1]$ is the additive identity, observe that

$$[a, b] + [0, 1] = [a1 + b0, b1] = [a, b].$$

It is easy to show that $[1, 1]$ is the multiplicative identity. Let $[a, b] \in F_D$ such that $a \neq 0$. Then $[b, a]$ is also in F_D and $[a, b] \cdot [b, a] = [1, 1]$; hence, $[b, a]$ is the multiplicative inverse for $[a, b]$. Similarly, $[-a, b]$ is the additive inverse of $[a, b]$. We leave as exercises the verification of the associative and commutative properties of multiplication in F_D. We also leave it to the reader to show that F_D is an abelian group under addition.

It remains to show that the distributive property holds in F_D; however,

$$
\begin{aligned}
[a, b][e, f] + [c, d][e, f] &= [ae, bf] + [ce, df] \\
&= [aedf + bfce, bdf^2] \\
&= [aed + bce, bdf] \\
&= [ade + bce, bdf] \\
&= ([a, b] + [c, d])[e, f]
\end{aligned}
$$

and the lemma is proved. $\square$

The field F_D in Lemma 18.3 is called the **field of fractions** or **field of quotients** of the integral domain D.

Theorem 18.4. *Let D be an integral domain. Then D can be embedded in a field of fractions F_D, where any element in F_D can be expressed as the quotient of two elements in D. Furthermore, the field of fractions F_D is unique in the sense that if E is any field containing D, then there exists a map $\psi : F_D \to E$ giving an isomorphism with a subfield of E such that $\psi(a) = a$ for all elements $a \in D$, where we identify a with its image in F_D.*

PROOF.   We will first demonstrate that D can be embedded in the field F_D. Define a map $\phi : D \to F_D$ by $\phi(a) = [a, 1]$. Then for a and b in D,

$$\phi(a + b) = [a + b, 1] = [a, 1] + [b, 1] = \phi(a) + \phi(b)$$

and

$$\phi(ab) = [ab, 1] = [a, 1][b, 1] = \phi(a)\phi(b);$$

hence, ϕ is a homomorphism. To show that ϕ is one-to-one, suppose that $\phi(a) = \phi(b)$. Then $[a, 1] = [b, 1]$, or $a = a1 = 1b = b$. Finally, any element of F_D can be expressed as the quotient of two elements in D, since

$$\phi(a)[\phi(b)]^{-1} = [a, 1][b, 1]^{-1} = [a, 1] \cdot [1, b] = [a, b].$$

Now let E be a field containing D and define a map $\psi : F_D \to E$ by $\psi([a, b]) = ab^{-1}$. To show that ψ is well-defined, let $[a_1, b_1] = [a_2, b_2]$. Then $a_1 b_2 = b_1 a_2$. Therefore, $a_1 b_1^{-1} = a_2 b_2^{-1}$ and $\psi([a_1, b_1]) = \psi([a_2, b_2])$.

If $[a, b]$ and $[c, d]$ are in F_D, then

$$
\begin{aligned}
\psi([a, b] + [c, d]) &= \psi([ad + bc, bd]) \\
&= (ad + bc)(bd)^{-1} \\
&= ab^{-1} + cd^{-1} \\
&= \psi([a, b]) + \psi([c, d])
\end{aligned}
$$

and

$$\psi([a,b] \cdot [c,d]) = \psi([ac,bd])$$
$$= (ac)(bd)^{-1}$$
$$= ab^{-1}cd^{-1}$$
$$= \psi([a,b])\psi([c,d]).$$

Therefore, ψ is a homomorphism.

To complete the proof of the theorem, we need to show that ψ is one-to-one. Suppose that $\psi([a,b]) = ab^{-1} = 0$. Then $a = 0b = 0$ and $[a,b] = [0,b]$. Therefore, the kernel of ψ is the zero element $[0,b]$ in F_D, and ψ is injective. $\qquad\square$

Example 18.5. Since $\mathbb{Q}$ is a field, $\mathbb{Q}[x]$ is an integral domain. The field of fractions of $\mathbb{Q}[x]$ is the set of all rational expressions $p(x)/q(x)$, where $p(x)$ and $q(x)$ are polynomials over the rationals and $q(x)$ is not the zero polynomial. We will denote this field by $\mathbb{Q}(x)$.

We will leave the proofs of the following corollaries of Theorem 18.4 as exercises.

Corollary 18.6. *Let F be a field of characteristic zero. Then F contains a subfield isomorphic to $\mathbb{Q}$.*

Corollary 18.7. *Let F be a field of characteristic p. Then F contains a subfield isomorphic to $\mathbb{Z}_p$.*

18.2 Factorization in Integral Domains

The building blocks of the integers are the prime numbers. If F is a field, then irreducible polynomials in $F[x]$ play a role that is very similar to that of the prime numbers in the ring of integers. Given an arbitrary integral domain, we are led to the following series of definitions.

Let R be a commutative ring with identity, and let a and b be elements in R. We say that a **divides** b, and write $a \mid b$, if there exists an element $c \in R$ such that $b = ac$. A **unit** in R is an element that has a multiplicative inverse. Two elements a and b in R are said to be **associates** if there exists a unit u in R such that $a = ub$.

Let D be an integral domain. A nonzero element $p \in D$ that is not a unit is said to be **irreducible** provided that whenever $p = ab$, either a or b is a unit. Furthermore, p is **prime** if whenever $p \mid ab$ either $p \mid a$ or $p \mid b$.

Example 18.8. It is important to notice that prime and irreducible elements do not always coincide. Let R be the subring (with identity) of $\mathbb{Q}[x,y]$ generated by x^2, y^2, and xy. Each of these elements is irreducible in R; however, xy is not prime, since xy divides x^2y^2 but does not divide either x^2 or y^2.

The Fundamental Theorem of Arithmetic states that every positive integer $n > 1$ can be factored into a product of prime numbers $p_1 \cdots p_k$, where the p_i's are not necessarily distinct. We also know that such factorizations are unique up to the order of the p_i's. We can easily extend this result to the integers. The question arises of whether or not such factorizations are possible in other rings. Generalizing this definition, we say an integral domain D is a **unique factorization domain**, or **ufd**, if D satisfies the following criteria.

1. Let $a \in D$ such that $a \neq 0$ and a is not a unit. Then a can be written as the product of irreducible elements in D.

2. Let $a = p_1 \cdots p_r = q_1 \cdots q_s$, where the p_i's and the q_i's are irreducible. Then $r = s$ and there is a $\pi \in S_r$ such that p_i and $q_{\pi(j)}$ are associates for $j = 1, \ldots, r$.

Example 18.9. The integers are a unique factorization domain by the Fundamental Theorem of Arithmetic.

Example 18.10. Not every integral domain is a unique factorization domain. The subring $\mathbb{Z}[\sqrt{3}\,i] = \{a + b\sqrt{3}\,i\}$ of the complex numbers is an integral domain (Exercise 16.6.12, Chapter 16). Let $z = a + b\sqrt{3}\,i$ and define $\nu : \mathbb{Z}[\sqrt{3}\,i] \to \mathbb{N} \cup \{0\}$ by $\nu(z) = |z|^2 = a^2 + 3b^2$. It is clear that $\nu(z) \geq 0$ with equality when $z = 0$. Also, from our knowledge of complex numbers we know that $\nu(zw) = \nu(z)\nu(w)$. It is easy to show that if $\nu(z) = 1$, then z is a unit, and that the only units of $\mathbb{Z}[\sqrt{3}\,i]$ are 1 and -1.

We claim that 4 has two distinct factorizations into irreducible elements:

$$4 = 2 \cdot 2 = (1 - \sqrt{3}\,i)(1 + \sqrt{3}\,i).$$

We must show that each of these factors is an irreducible element in $\mathbb{Z}[\sqrt{3}\,i]$. If 2 is not irreducible, then $2 = zw$ for elements z, w in $\mathbb{Z}[\sqrt{3}\,i]$ where $\nu(z) = \nu(w) = 2$. However, there does not exist an element in z in $\mathbb{Z}[\sqrt{3}\,i]$ such that $\nu(z) = 2$ because the equation $a^2 + 3b^2 = 2$ has no integer solutions. Therefore, 2 must be irreducible. A similar argument shows that both $1 - \sqrt{3}\,i$ and $1 + \sqrt{3}\,i$ are irreducible. Since 2 is not a unit multiple of either $1 - \sqrt{3}\,i$ or $1 + \sqrt{3}\,i$, 4 has at least two distinct factorizations into irreducible elements.

Principal Ideal Domains

Let R be a commutative ring with identity. Recall that a principal ideal generated by $a \in R$ is an ideal of the form $\langle a \rangle = \{ra : r \in R\}$. An integral domain in which every ideal is principal is called a ***principal ideal domain***, or **pid**.

Lemma 18.11. *Let D be an integral domain and let $a, b \in D$. Then*

1. *$a \mid b$ if and only if $\langle b \rangle \subset \langle a \rangle$.*

2. *a and b are associates if and only if $\langle b \rangle = \langle a \rangle$.*

3. *a is a unit in D if and only if $\langle a \rangle = D$.*

PROOF.   (1) Suppose that $a \mid b$. Then $b = ax$ for some $x \in D$. Hence, for every r in D, $br = (ax)r = a(xr)$ and $\langle b \rangle \subset \langle a \rangle$. Conversely, suppose that $\langle b \rangle \subset \langle a \rangle$. Then $b \in \langle a \rangle$. Consequently, $b = ax$ for some $x \in D$. Thus, $a \mid b$.

(2) Since a and b are associates, there exists a unit u such that $a = ub$. Therefore, $b \mid a$ and $\langle a \rangle \subset \langle b \rangle$. Similarly, $\langle b \rangle \subset \langle a \rangle$. It follows that $\langle a \rangle = \langle b \rangle$. Conversely, suppose that $\langle a \rangle = \langle b \rangle$. By part (1), $a \mid b$ and $b \mid a$. Then $a = bx$ and $b = ay$ for some $x, y \in D$. Therefore, $a = bx = ayx$. Since D is an integral domain, $xy = 1$; that is, x and y are units and a and b are associates.

(3) An element $a \in D$ is a unit if and only if a is an associate of 1. However, a is an associate of 1 if and only if $\langle a \rangle = \langle 1 \rangle = D$. $\square$

Theorem 18.12. *Let D be a* PID *and $\langle p \rangle$ be a nonzero ideal in D. Then $\langle p \rangle$ is a maximal ideal if and only if p is irreducible.*

PROOF.   Suppose that $\langle p \rangle$ is a maximal ideal. If some element a in D divides p, then $\langle p \rangle \subset \langle a \rangle$. Since $\langle p \rangle$ is maximal, either $D = \langle a \rangle$ or $\langle p \rangle = \langle a \rangle$. Consequently, either a and p are associates or a is a unit. Therefore, p is irreducible.

Conversely, let p be irreducible. If $\langle a \rangle$ is an ideal in D such that $\langle p \rangle \subset \langle a \rangle \subset D$, then $a \mid p$. Since p is irreducible, either a must be a unit or a and p are associates. Therefore, either $D = \langle a \rangle$ or $\langle p \rangle = \langle a \rangle$. Thus, $\langle p \rangle$ is a maximal ideal. $\square$

Corollary 18.13. *Let D be a* PID. *If p is irreducible, then p is prime.*

PROOF. Let p be irreducible and suppose that $p \mid ab$. Then $\langle ab \rangle \subset \langle p \rangle$. By Corollary 16.40, since $\langle p \rangle$ is a maximal ideal, $\langle p \rangle$ must also be a prime ideal. Thus, either $a \in \langle p \rangle$ or $b \in \langle p \rangle$. Hence, either $p \mid a$ or $p \mid b$. $\qquad \square$

Lemma 18.14. *Let D be a* PID. *Let $I_1, I_2, \ldots$ be a set of ideals such that $I_1 \subset I_2 \subset \cdots$. Then there exists an integer N such that $I_n = I_N$ for all $n \geq N$.*

PROOF. We claim that $I = \bigcup_{i=1}^{\infty} I_i$ is an ideal of D. Certainly I is not empty, since $I_1 \subset I$ and $0 \in I$. If $a, b \in I$, then $a \in I_i$ and $b \in I_j$ for some i and j in $\mathbb{N}$. Without loss of generality we can assume that $i \leq j$. Hence, a and b are both in I_j and so $a - b$ is also in I_j. Now let $r \in D$ and $a \in I$. Again, we note that $a \in I_i$ for some positive integer i. Since I_i is an ideal, $ra \in I_i$ and hence must be in I. Therefore, we have shown that I is an ideal in D.

Since D is a principal ideal domain, there exists an element $\bar{a} \in D$ that generates I. Since $\bar{a}$ is in I_N for some $N \in \mathbb{N}$, we know that $I_N = I = \langle \bar{a} \rangle$. Consequently, $I_n = I_N$ for $n \geq N$. $\qquad \square$

Any commutative ring satisfying the condition in Lemma 18.14 is said to satisfy the **ascending chain condition**, or **ACC**. Such rings are called **Noetherian rings**, after Emmy Noether.

Theorem 18.15. *Every* PID *is a* UFD.

PROOF. *Existence of a factorization.* Let D be a PID and a be a nonzero element in D that is not a unit. If a is irreducible, then we are done. If not, then there exists a factorization $a = a_1 b_1$, where neither a_1 nor b_1 is a unit. Hence, $\langle a \rangle \subset \langle a_1 \rangle$. By Lemma 18.11, we know that $\langle a \rangle \neq \langle a_1 \rangle$; otherwise, a and a_1 would be associates and b_1 would be a unit, which would contradict our assumption. Now suppose that $a_1 = a_2 b_2$, where neither a_2 nor b_2 is a unit. By the same argument as before, $\langle a_1 \rangle \subset \langle a_2 \rangle$. We can continue with this construction to obtain an ascending chain of ideals

$$\langle a \rangle \subset \langle a_1 \rangle \subset \langle a_2 \rangle \subset \cdots .$$

By Lemma 18.14, there exists a positive integer N such that $\langle a_n \rangle = \langle a_N \rangle$ for all $n \geq N$. Consequently, a_N must be irreducible. We have now shown that a is the product of two elements, one of which must be irreducible.

Now suppose that $a = c_1 p_1$, where p_1 is irreducible. If c_1 is not a unit, we can repeat the preceding argument to conclude that $\langle a \rangle \subset \langle c_1 \rangle$. Either c_1 is irreducible or $c_1 = c_2 p_2$, where p_2 is irreducible and c_2 is not a unit. Continuing in this manner, we obtain another chain of ideals

$$\langle a \rangle \subset \langle c_1 \rangle \subset \langle c_2 \rangle \subset \cdots .$$

This chain must satisfy the ascending chain condition; therefore,

$$a = p_1 p_2 \cdots p_r$$

for irreducible elements $p_1, \ldots, p_r$.

Uniqueness of the factorization. To show uniqueness, let

$$a = p_1 p_2 \cdots p_r = q_1 q_2 \cdots q_s,$$

where each p_i and each q_i is irreducible. Without loss of generality, we can assume that $r < s$. Since p_1 divides $q_1 q_2 \cdots q_s$, by Corollary 18.13 it must divide some q_i. By rearranging the q_i's, we can assume that $p_1 \mid q_1$; hence, $q_1 = u_1 p_1$ for some unit u_1 in D. Therefore,

$$a = p_1 p_2 \cdots p_r = u_1 p_1 q_2 \cdots q_s$$

or

$$p_2 \cdots p_r = u_1 q_2 \cdots q_s.$$

Continuing in this manner, we can arrange the q_i's such that $p_2 = q_2, p_3 = q_3, \ldots, p_r = q_r$, to obtain

$$u_1 u_2 \cdots u_r q_{r+1} \cdots q_s = 1.$$

In this case $q_{r+1} \cdots q_s$ is a unit, which contradicts the fact that $q_{r+1}, \ldots, q_s$ are irreducibles. Therefore, $r = s$ and the factorization of a is unique. $\qquad\square$

Corollary 18.16. *Let F be a field. Then $F[x]$ is a* UFD*.*

Example 18.17. Every PID is a UFD, but it is not the case that every UFD is a PID. In Corollary 18.31, we will prove that $\mathbb{Z}[x]$ is a UFD. However, $\mathbb{Z}[x]$ is not a PID. Let $I = \{5f(x) + xg(x) : f(x), g(x) \in \mathbb{Z}[x]\}$. We can easily show that I is an ideal of $\mathbb{Z}[x]$. Suppose that $I = \langle p(x) \rangle$. Since $5 \in I$, $5 = f(x)p(x)$. In this case $p(x) = p$ must be a constant. Since $x \in I$, $x = pg(x)$; consequently, $p = \pm 1$. However, it follows from this fact that $\langle p(x) \rangle = \mathbb{Z}[x]$. But this would mean that 3 is in I. Therefore, we can write $3 = 5f(x) + xg(x)$ for some $f(x)$ and $g(x)$ in $\mathbb{Z}[x]$. Examining the constant term of this polynomial, we see that $3 = 5f(x)$, which is impossible.

Euclidean Domains

We have repeatedly used the division algorithm when proving results about either $\mathbb{Z}$ or $F[x]$, where F is a field. We should now ask when a division algorithm is available for an integral domain.

Let D be an integral domain such that for each $a \in D$ there is a nonnegative integer $\nu(a)$ satisfying the following conditions.

1. If a and b are nonzero elements in D, then $\nu(a) \leq \nu(ab)$.

2. Let $a, b \in D$ and suppose that $b \neq 0$. Then there exist elements $q, r \in D$ such that $a = bq + r$ and either $r = 0$ or $\nu(r) < \nu(b)$.

Then D is called a **Euclidean domain** and ν is called a **Euclidean valuation**.

Example 18.18. Absolute value on $\mathbb{Z}$ is a Euclidean valuation.

Example 18.19. Let F be a field. Then the degree of a polynomial in $F[x]$ is a Euclidean valuation.

Example 18.20. Recall that the Gaussian integers in Example 16.12 of Chapter 16 are defined by

$$\mathbb{Z}[i] = \{a + bi : a, b \in \mathbb{Z}\}.$$

We usually measure the size of a complex number $a + bi$ by its absolute value, $|a + bi| = \sqrt{a^2 + b^2}$; however, $\sqrt{a^2 + b^2}$ may not be an integer. For our valuation we will let $\nu(a+bi) = a^2 + b^2$ to ensure that we have an integer.

We claim that $\nu(a + bi) = a^2 + b^2$ is a Euclidean valuation on $\mathbb{Z}[i]$. Let $z, w \in \mathbb{Z}[i]$. Then $\nu(zw) = |zw|^2 = |z|^2|w|^2 = \nu(z)\nu(w)$. Since $\nu(z) \geq 1$ for every nonzero $z \in \mathbb{Z}[i]$, $\nu(z) \leq \nu(z)\nu(w)$.

Next, we must show that for any $z = a + bi$ and $w = c + di$ in $\mathbb{Z}[i]$ with $w \neq 0$, there exist elements q and r in $\mathbb{Z}[i]$ such that $z = qw + r$ with either $r = 0$ or $\nu(r) < \nu(w)$. We

can view z and w as elements in $\mathbb{Q}(i) = \{p + qi : p, q \in \mathbb{Q}\}$, the field of fractions of $\mathbb{Z}[i]$. Observe that

$$\begin{aligned}
zw^{-1} &= (a + bi)\frac{c - di}{c^2 + d^2} \\
&= \frac{ac + bd}{c^2 + d^2} + \frac{bc - ad}{c^2 + d^2}i \\
&= \left(m_1 + \frac{n_1}{c^2 + d^2}\right) + \left(m_2 + \frac{n_2}{c^2 + d^2}\right)i \\
&= (m_1 + m_2 i) + \left(\frac{n_1}{c^2 + d^2} + \frac{n_2}{c^2 + d^2}i\right) \\
&= (m_1 + m_2 i) + (s + ti)
\end{aligned}$$

in $\mathbb{Q}(i)$. In the last steps we are writing the real and imaginary parts as an integer plus a proper fraction. That is, we take the closest integer m_i such that the fractional part satisfies $|n_i/(a^2 + b^2)| \leq 1/2$. For example, we write

$$\begin{aligned}
\frac{9}{8} &= 1 + \frac{1}{8} \\
\frac{15}{8} &= 2 - \frac{1}{8}.
\end{aligned}$$

Thus, s and t are the "fractional parts" of $zw^{-1} = (m_1 + m_2 i) + (s + ti)$. We also know that $s^2 + t^2 \leq 1/4 + 1/4 = 1/2$. Multiplying by w, we have

$$z = zw^{-1}w = w(m_1 + m_2 i) + w(s + ti) = qw + r,$$

where $q = m_1 + m_2 i$ and $r = w(s + ti)$. Since z and qw are in $\mathbb{Z}[i]$, r must be in $\mathbb{Z}[i]$. Finally, we need to show that either $r = 0$ or $\nu(r) < \nu(w)$. However,

$$\nu(r) = \nu(w)\nu(s + ti) \leq \frac{1}{2}\nu(w) < \nu(w).$$

Theorem 18.21. *Every Euclidean domain is a principal ideal domain.*

PROOF. Let D be a Euclidean domain and let ν be a Euclidean valuation on D. Suppose I is a nontrivial ideal in D and choose a nonzero element $b \in I$ such that $\nu(b)$ is minimal for all $a \in I$. Since D is a Euclidean domain, there exist elements q and r in D such that $a = bq + r$ and either $r = 0$ or $\nu(r) < \nu(b)$. But $r = a - bq$ is in I since I is an ideal; therefore, $r = 0$ by the minimality of b. It follows that $a = bq$ and $I = \langle b \rangle$. $\square$

Corollary 18.22. *Every Euclidean domain is a unique factorization domain.*

Factorization in $D[x]$

One of the most important polynomial rings is $\mathbb{Z}[x]$. One of the first questions that come to mind about $\mathbb{Z}[x]$ is whether or not it is a UFD. We will prove a more general statement here. Our first task is to obtain a more general version of Gauss's Lemma (Theorem 17.14).

Let D be a unique factorization domain and suppose that

$$p(x) = a_n x^n + \cdots + a_1 x + a_0$$

in $D[x]$. Then the **content** of $p(x)$ is the greatest common divisor of $a_0, \ldots, a_n$. We say that $p(x)$ is **primitive** if $\gcd(a_0, \ldots, a_n) = 1$.

Example 18.23. In $\mathbb{Z}[x]$ the polynomial $p(x) = 5x^4 - 3x^3 + x - 4$ is a primitive polynomial since the greatest common divisor of the coefficients is 1; however, the polynomial $q(x) = 4x^2 - 6x + 8$ is not primitive since the content of $q(x)$ is 2.

Theorem 18.24 Gauss's Lemma. *Let D be a* UFD *and let $f(x)$ and $g(x)$ be primitive polynomials in $D[x]$. Then $f(x)g(x)$ is primitive.*

PROOF.   Let $f(x) = \sum_{i=0}^{m} a_i x^i$ and $g(x) = \sum_{i=0}^{n} b_i x^i$. Suppose that p is a prime dividing the coefficients of $f(x)g(x)$. Let r be the smallest integer such that $p \nmid a_r$ and s be the smallest integer such that $p \nmid b_s$. The coefficient of x^{r+s} in $f(x)g(x)$ is

$$c_{r+s} = a_0 b_{r+s} + a_1 b_{r+s-1} + \cdots + a_{r+s-1} b_1 + a_{r+s} b_0.$$

Since p divides $a_0, \ldots, a_{r-1}$ and $b_0, \ldots, b_{s-1}$, p divides every term of c_{r+s} except for the term $a_r b_s$. However, since $p \mid c_{r+s}$, either p divides a_r or p divides b_s. But this is impossible. $\square$

Lemma 18.25. *Let D be a* UFD, *and let $p(x)$ and $q(x)$ be in $D[x]$. Then the content of $p(x)q(x)$ is equal to the product of the contents of $p(x)$ and $q(x)$.*

PROOF.   Let $p(x) = cp_1(x)$ and $q(x) = dq_1(x)$, where c and d are the contents of $p(x)$ and $q(x)$, respectively. Then $p_1(x)$ and $q_1(x)$ are primitive. We can now write $p(x)q(x) = cdp_1(x)q_1(x)$. Since $p_1(x)q_1(x)$ is primitive, the content of $p(x)q(x)$ must be cd. $\square$

Lemma 18.26. *Let D be a* UFD *and F its field of fractions. Suppose that $p(x) \in D[x]$ and $p(x) = f(x)g(x)$, where $f(x)$ and $g(x)$ are in $F[x]$. Then $p(x) = f_1(x)g_1(x)$, where $f_1(x)$ and $g_1(x)$ are in $D[x]$. Furthermore, $\deg f(x) = \deg f_1(x)$ and $\deg g(x) = \deg g_1(x)$.*

PROOF.   Let a and b be nonzero elements of D such that $af(x), bg(x)$ are in $D[x]$. We can find $a_1, b_2 \in D$ such that $af(x) = a_1 f_1(x)$ and $bg(x) = b_1 g_1(x)$, where $f_1(x)$ and $g_1(x)$ are primitive polynomials in $D[x]$. Therefore, $abp(x) = (a_1 f_1(x))(b_1 g_1(x))$. Since $f_1(x)$ and $g_1(x)$ are primitive polynomials, it must be the case that $ab \mid a_1 b_1$ by Gauss's Lemma. Thus there exists a $c \in D$ such that $p(x) = cf_1(x)g_1(x)$. Clearly, $\deg f(x) = \deg f_1(x)$ and $\deg g(x) = \deg g_1(x)$. $\square$

The following corollaries are direct consequences of Lemma 18.26.

Corollary 18.27. *Let D be a* UFD *and F its field of fractions. A primitive polynomial $p(x)$ in $D[x]$ is irreducible in $F[x]$ if and only if it is irreducible in $D[x]$.*

Corollary 18.28. *Let D be a* UFD *and F its field of fractions. If $p(x)$ is a monic polynomial in $D[x]$ with $p(x) = f(x)g(x)$ in $F[x]$, then $p(x) = f_1(x)g_1(x)$, where $f_1(x)$ and $g_1(x)$ are in $D[x]$. Furthermore, $\deg f(x) = \deg f_1(x)$ and $\deg g(x) = \deg g_1(x)$.*

Theorem 18.29. *If D is a* UFD, *then $D[x]$ is a* UFD.

PROOF.   Let $p(x)$ be a nonzero polynomial in $D[x]$. If $p(x)$ is a constant polynomial, then it must have a unique factorization since D is a UFD. Now suppose that $p(x)$ is a polynomial of positive degree in $D[x]$. Let F be the field of fractions of D, and let $p(x) = f_1(x) f_2(x) \cdots f_n(x)$ by a factorization of $p(x)$, where each $f_i(x)$ is irreducible. Choose $a_i \in D$ such that $a_i f_i(x)$ is in $D[x]$. There exist $b_1, \ldots, b_n \in D$ such that $a_i f_i(x) = b_i g_i(x)$, where $g_i(x)$ is a primitive polynomial in $D[x]$. By Corollary 18.27, each $g_i(x)$ is irreducible in $D[x]$. Consequently, we can write

$$a_1 \cdots a_n p(x) = b_1 \cdots b_n g_1(x) \cdots g_n(x).$$

Let $b = b_1 \cdots b_n$. Since $g_1(x) \cdots g_n(x)$ is primitive, $a_1 \cdots a_n$ divides b. Therefore, $p(x) = ag_1(x) \cdots g_n(x)$, where $a \in D$. Since D is a UFD, we can factor a as $uc_1 \cdots c_k$, where u is a unit and each of the c_i's is irreducible in D.

We will now show the uniqueness of this factorization. Let

$$p(x) = a_1 \cdots a_m f_1(x) \cdots f_n(x) = b_1 \cdots b_r g_1(x) \cdots g_s(x)$$

be two factorizations of $p(x)$, where all of the factors are irreducible in $D[x]$. By Corollary 18.27, each of the f_i's and g_i's is irreducible in $F[x]$. The a_i's and the b_i's are units in F. Since $F[x]$ is a PID, it is a UFD; therefore, $n = s$. Now rearrange the $g_i(x)$'s so that $f_i(x)$ and $g_i(x)$ are associates for $i = 1, \ldots, n$. Then there exist $c_1, \ldots, c_n$ and $d_1, \ldots, d_n$ in D such that $(c_i/d_i)f_i(x) = g_i(x)$ or $c_i f_i(x) = d_i g_i(x)$. The polynomials $f_i(x)$ and $g_i(x)$ are primitive; hence, c_i and d_i are associates in D. Thus, $a_1 \cdots a_m = ub_1 \cdots b_r$ in D, where u is a unit in D. Since D is a unique factorization domain, $m = s$. Finally, we can reorder the b_i's so that a_i and b_i are associates for each i. This completes the uniqueness part of the proof. $\square$

The theorem that we have just proven has several obvious but important corollaries.

Corollary 18.30. *Let F be a field. Then $F[x]$ is a UFD.*

Corollary 18.31. *The ring of polynomials over the integers, $\mathbb{Z}[x]$, is a UFD.*

Corollary 18.32. *Let D be a UFD. Then $D[x_1, \ldots, x_n]$ is a UFD.*

Remark 18.33. It is important to notice that every Euclidean domain is a PID and every PID is a UFD. However, as demonstrated by our examples, the converse of each of these statements fails. There are principal ideal domains that are not Euclidean domains, and there are unique factorization domains that are not principal ideal domains ($\mathbb{Z}[x]$).

Historical Note

Karl Friedrich Gauss, born in Brunswick, Germany on April 30, 1777, is considered to be one of the greatest mathematicians who ever lived. Gauss was truly a child prodigy. At the age of three he was able to detect errors in the books of his father's business. Gauss entered college at the age of 15. Before the age of 20, Gauss was able to construct a regular 17-sided polygon with a ruler and compass. This was the first new construction of a regular n-sided polygon since the time of the ancient Greeks. Gauss succeeded in showing that if $N = 2^{2^n} + 1$ was prime, then it was possible to construct a regular N-sided polygon.

Gauss obtained his Ph.D. in 1799 under the direction of Pfaff at the University of Helmstedt. In his dissertation he gave the first complete proof of the Fundamental Theorem of Algebra, which states that every polynomial with real coefficients can be factored into linear factors over the complex numbers. The acceptance of complex numbers was brought about by Gauss, who was the first person to use the notation of i for $\sqrt{-1}$.

Gauss then turned his attention toward number theory; in 1801, he published his famous book on number theory, *Disquisitiones Arithmeticae*. Throughout his life Gauss was intrigued with this branch of mathematics. He once wrote, "Mathematics is the queen of the sciences, and the theory of numbers is the queen of mathematics."

In 1807, Gauss was appointed director of the Observatory at the University of Göttingen, a position he held until his death. This position required him to study applications of mathematics to the sciences. He succeeded in making contributions to fields such as astronomy, mechanics, optics, geodesy, and magnetism. Along with Wilhelm Weber, he coinvented the

first practical electric telegraph some years before a better version was invented by Samuel F. B. Morse.

Gauss was clearly the most prominent mathematician in the world in the early nineteenth century. His status naturally made his discoveries subject to intense scrutiny. Gauss's cold and distant personality many times led him to ignore the work of his contemporaries, making him many enemies. He did not enjoy teaching very much, and young mathematicians who sought him out for encouragement were often rebuffed. Nevertheless, he had many outstanding students, including Eisenstein, Riemann, Kummer, Dirichlet, and Dedekind. Gauss also offered a great deal of encouragement to Sophie Germain (1776–1831), who overcame the many obstacles facing women in her day to become a very prominent mathematician. Gauss died at the age of 78 in Göttingen on February 23, 1855.

18.3   Exercises

1.   Let $z = a + b\sqrt{3}\,i$ be in $\mathbb{Z}[\sqrt{3}\,i]$. If $a^2 + 3b^2 = 1$, show that z must be a unit. Show that the only units of $\mathbb{Z}[\sqrt{3}\,i]$ are 1 and -1.

2.   The Gaussian integers, $\mathbb{Z}[i]$, are a UFD. Factor each of the following elements in $\mathbb{Z}[i]$ into a product of irreducibles.

  (a) 5

  (b) $1 + 3i$

  (c) $6 + 8i$

  (d) 2

3.   Let D be an integral domain.

  (a) Prove that F_D is an abelian group under the operation of addition.

  (b) Show that the operation of multiplication is well-defined in the field of fractions, F_D.

  (c) Verify the associative and commutative properties for multiplication in F_D.

4.   Prove or disprove: Any subring of a field F containing 1 is an integral domain.

5.   Prove or disprove: If D is an integral domain, then every prime element in D is also irreducible in D.

6.   Let F be a field of characteristic zero. Prove that F contains a subfield isomorphic to $\mathbb{Q}$.

7.   Let F be a field.

  (a) Prove that the field of fractions of $F[x]$, denoted by $F(x)$, is isomorphic to the set all rational expressions $p(x)/q(x)$, where $q(x)$ is not the zero polynomial.

  (b) Let $p(x_1, \ldots, x_n)$ and $q(x_1, \ldots, x_n)$ be polynomials in $F[x_1, \ldots, x_n]$. Show that the set of all rational expressions $p(x_1, \ldots, x_n)/q(x_1, \ldots, x_n)$ is isomorphic to the field of fractions of $F[x_1, \ldots, x_n]$. We denote the field of fractions of $F[x_1, \ldots, x_n]$ by $F(x_1, \ldots, x_n)$.

8.   Let p be prime and denote the field of fractions of $\mathbb{Z}_p[x]$ by $\mathbb{Z}_p(x)$. Prove that $\mathbb{Z}_p(x)$ is an infinite field of characteristic p.

9.   Prove that the field of fractions of the Gaussian integers, $\mathbb{Z}[i]$, is

$$\mathbb{Q}(i) = \{p + qi : p, q \in \mathbb{Q}\}.$$

10.   A field F is called a ***prime field*** if it has no proper subfields. If E is a subfield of F and E is a prime field, then E is a ***prime subfield*** of F.

(a) Prove that every field contains a unique prime subfield.

(b) If F is a field of characteristic 0, prove that the prime subfield of F is isomorphic to the field of rational numbers, $\mathbb{Q}$.

(c) If F is a field of characteristic p, prove that the prime subfield of F is isomorphic to $\mathbb{Z}_p$.

11.   Let $\mathbb{Z}[\sqrt{2}\,] = \{a + b\sqrt{2} : a, b \in \mathbb{Z}\}$.

(a) Prove that $\mathbb{Z}[\sqrt{2}\,]$ is an integral domain.

(b) Find all of the units in $\mathbb{Z}[\sqrt{2}\,]$.

(c) Determine the field of fractions of $\mathbb{Z}[\sqrt{2}\,]$.

(d) Prove that $\mathbb{Z}[\sqrt{2}i]$ is a Euclidean domain under the Euclidean valuation $\nu(a + b\sqrt{2}\,i) = a^2 + 2b^2$.

12.   Let D be a UFD. An element $d \in D$ is a **greatest common divisor of** a **and** b **in** D if $d \mid a$ and $d \mid b$ and d is divisible by any other element dividing both a and b.

(a) If D is a PID and a and b are both nonzero elements of D, prove there exists a unique greatest common divisor of a and b up to associates. That is, if d and d' are both greatest common divisors of a and b, then d and d' are associates. We write $\gcd(a, b)$ for the greatest common divisor of a and b.

(b) Let D be a PID and a and b be nonzero elements of D. Prove that there exist elements s and t in D such that $\gcd(a, b) = as + bt$.

13.   Let D be an integral domain. Define a relation on D by $a \sim b$ if a and b are associates in D. Prove that $\sim$ is an equivalence relation on D.

14.   Let D be a Euclidean domain with Euclidean valuation ν. If u is a unit in D, show that $\nu(u) = \nu(1)$.

15.   Let D be a Euclidean domain with Euclidean valuation ν. If a and b are associates in D, prove that $\nu(a) = \nu(b)$.

16.   Show that $\mathbb{Z}[\sqrt{5}\,i]$ is not a unique factorization domain.

17.   Prove or disprove: Every subdomain of a UFD is also a UFD.

18.   An ideal of a commutative ring R is said to be **finitely generated** if there exist elements $a_1, \ldots, a_n$ in R such that every element $r \in R$ can be written as $a_1 r_1 + \cdots + a_n r_n$ for some $r_1, \ldots, r_n$ in R. Prove that R satisfies the ascending chain condition if and only if every ideal of R is finitely generated.

19.   Let D be an integral domain with a descending chain of ideals $I_1 \supset I_2 \supset I_3 \supset \cdots$. Suppose that there exists an N such that $I_k = I_N$ for all $k \geq N$. A ring satisfying this condition is said to satisfy the **descending chain condition**, or **DCC**. Rings satisfying the DCC are called **Artinian rings**, after Emil Artin. Show that if D satisfies the descending chain condition, it must satisfy the ascending chain condition.

20.   Let R be a commutative ring with identity. We define a **multiplicative subset** of R to be a subset S such that $1 \in S$ and $ab \in S$ if $a, b \in S$.

(a) Define a relation $\sim$ on $R \times S$ by $(a, s) \sim (a', s')$ if there exists an $s^* \in S$ such that $s^*(s'a - sa') = 0$. Show that $\sim$ is an equivalence relation on $R \times S$.

(b) Let a/s denote the equivalence class of $(a, s) \in R \times S$ and let $S^{-1}R$ be the set of all equivalence classes with respect to $\sim$. Define the operations of addition and multiplication on $S^{-1}R$ by

$$\frac{a}{s} + \frac{b}{t} = \frac{at + bs}{st}$$
$$\frac{a}{s}\frac{b}{t} = \frac{ab}{st},$$

respectively. Prove that these operations are well-defined on $S^{-1}R$ and that $S^{-1}R$ is a ring with identity under these operations. The ring $S^{-1}R$ is called the ***ring of quotients*** of R with respect to S.

(c) Show that the map $\psi : R \rightarrow S^{-1}R$ defined by $\psi(a) = a/1$ is a ring homomorphism.

(d) If R has no zero divisors and $0 \notin S$, show that ψ is one-to-one.

(e) Prove that P is a prime ideal of R if and only if $S = R \setminus P$ is a multiplicative subset of R.

(f) If P is a prime ideal of R and $S = R \setminus P$, show that the ring of quotients $S^{-1}R$ has a unique maximal ideal. Any ring that has a unique maximal ideal is called a ***local ring***.

References and Suggested Readings

[**1**]   Atiyah, M. F. and MacDonald, I. G. *Introduction to Commutative Algebra*. Westview Press, Boulder, CO, 1994.

[**2**]   Zariski, O. and Samuel, P. *Commutative Algebra*, vols. I and II. Springer, New York, 1975, 1960.

18.4   Sage

We have already seen some integral domains and unique factorizations in the previous two chapters. In addition to what we have already seen, Sage has support for some of the topics from this section, but the coverage is limited. Some functions will work for some rings and not others, while some functions are not yet part of Sage. So we will give some examples, but this is far from comprehensive.

Field of Fractions

Sage is frequently able to construct a field of fractions, or identify a certain field as the field of fractions. For example, the ring of integers and the field of rational numbers are both implemented in Sage, and the integers "know" that the rationals is it's field of fractions.

```
Q = ZZ.fraction_field(); Q
```

```
Rational Field
```

```
Q == QQ
```

True

In other cases Sage will construct a fraction field, in the spirit of Lemma 18.3. So it is then possible to do basic calculations in the constructed field.

```
R.<x> = ZZ[]
P = R.fraction_field();P
```

Fraction Field of Univariate Polynomial Ring **in** x over Integer Ring

```
f = P((x^2+3)/(7*x+4))
g = P((4*x^2)/(3*x^2-5*x+4))
h = P((-2*x^3+4*x^2+3)/(x^2+1))
((f+g)/h).numerator()
```

3*x^6 + 23*x^5 + 32*x^4 + 8*x^3 + 41*x^2 - 15*x + 12

```
((f+g)/h).denominator()
```

-42*x^6 + 130*x^5 - 108*x^4 + 63*x^3 - 5*x^2 + 24*x + 48

Prime Subfields

Corollary 18.7 says every field of characteristic p has a subfield isomorphic to $\mathbb{Z}_p$. For a finite field, the exact nature of this subfield is not a surprise, but Sage will allow us to extract it easily.

```
F.<c> = FiniteField(3^5)
F.characteristic()
```

3

```
G = F.prime_subfield(); G
```

Finite Field of size 3

```
G.list()
```

[0, 1, 2]

More generally, the fields mentioned in the conclusions of Corollary 18.6 and Corollary 18.7 are known as the "prime subfield" of the ring containing them. Here is an example of the characteristic zero case.

```
K.<y>=QuadraticField(-7); K
```

Number Field **in** y with defining polynomial x^2 + 7

```
K.prime_subfield()
```

Rational Field

In a rough sense, every characteristic zero field contains a copy of the rational numbers (the fraction field of the integers), which can explain Sage's extensive support for rings and fields that extend the integers and the rationals.

Integral Domains

Sage can determine if some rings are integral domains and we can test products in them. However, notions of units, irreducibles or prime elements are not generally supported (outside of what we have seen for polynomials in the previous chapter). Worse, the construction below creates a ring within a larger field and so some functions (such as `.is_unit()`) pass through and give misleading results. This is because the construction below creates a ring known as an "order in a number field."

```
K.<x> = ZZ[sqrt(-3)]; K
```

```
Order in Number Field in a with defining polynomial x^2 + 3
```

```
K.is_integral_domain()
```

```
True
```

```
K.basis()
```

```
[1, a]
```

```
x
```

```
a
```

```
(1+x)*(1-x) == 2*2
```

```
True
```

The following is a bit misleading, since 4, as an element of $\mathbb{Z}[\sqrt{3}i]$ does not have a multiplicative inverse, though seemingly we can compute one.

```
four = K(4)
four.is_unit()
```

```
False
```

```
four^-1
```

```
1/4
```

Principal Ideals

When a ring is a principal ideal domain, such as the integers, or polynomials over a field, Sage works well. Beyond that, support begins to weaken.

```
T.<x>=ZZ[]
T.is_integral_domain()
```

```
True
```

```
J = T.ideal(5, x); J
```

```
Ideal (5, x) of Univariate Polynomial Ring in x over Integer Ring
```

```
Q = T.quotient(J); Q
```

```
Quotient of Univariate Polynomial Ring in x over
Integer Ring by the ideal (5, x)
```

```
J.is_principal()
```

```
Traceback (most recent call last):
...
NotImplementedError
```

```
Q.is_field()
```

```
Traceback (most recent call last):
...
NotImplementedError
```

18.5   Sage Exercises

There are no Sage exercises for this section.

19

Lattices and Boolean Algebras

The axioms of a ring give structure to the operations of addition and multiplication on a set. However, we can construct algebraic structures, known as lattices and Boolean algebras, that generalize other types of operations. For example, the important operations on sets are inclusion, union, and intersection. Lattices are generalizations of order relations on algebraic spaces, such as set inclusion in set theory and inequality in the familiar number systems $\mathbb{N}$, $\mathbb{Z}$, $\mathbb{Q}$, and $\mathbb{R}$. Boolean algebras generalize the operations of intersection and union. Lattices and Boolean algebras have found applications in logic, circuit theory, and probability.

19.1 Lattices

Partially Ordered Sets

We begin the study of lattices and Boolean algebras by generalizing the idea of inequality. Recall that a **relation** on a set X is a subset of $X \times X$. A relation P on X is called a **partial order** of X if it satisfies the following axioms.

1. The relation is **reflexive**: $(a, a) \in P$ for all $a \in X$.

2. The relation is **antisymmetric**: if $(a, b) \in P$ and $(b, a) \in P$, then $a = b$.

3. The relation is **transitive**: if $(a, b) \in P$ and $(b, c) \in P$, then $(a, c) \in P$.

We will usually write $a \preceq b$ to mean $(a, b) \in P$ unless some symbol is naturally associated with a particular partial order, such as $a \leq b$ with integers a and b, or $A \subset B$ with sets A and B. A set X together with a partial order $\preceq$ is called a **partially ordered set**, or **poset**.

Example 19.1. The set of integers (or rationals or reals) is a poset where $a \leq b$ has the usual meaning for two integers a and b in $\mathbb{Z}$.

Example 19.2. Let X be any set. We will define the **power set** of X to be the set of all subsets of X. We denote the power set of X by $\mathcal{P}(X)$. For example, let $X = \{a, b, c\}$. Then $\mathcal{P}(X)$ is the set of all subsets of the set $\{a, b, c\}$:

$$\emptyset \qquad \{a\} \qquad \{b\} \qquad \{c\}$$
$$\{a, b\} \qquad \{a, c\} \qquad \{b, c\} \qquad \{a, b, c\}.$$

On any power set of a set X, set inclusion, $\subset$, is a partial order. We can represent the order on $\{a, b, c\}$ schematically by a diagram such as the one in Figure 19.3.
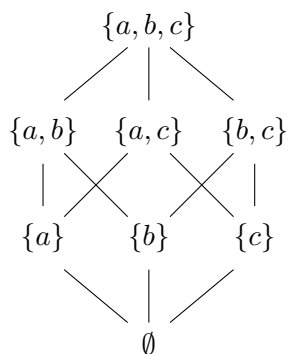
Figure 19.3: Partial order on $\mathcal{P}(\{a, b, c\})$

Example 19.4. Let G be a group. The set of subgroups of G is a poset, where the partial order is set inclusion.

Example 19.5. There can be more than one partial order on a particular set. We can form a partial order on $\mathbb{N}$ by $a \preceq b$ if $a \mid b$. The relation is certainly reflexive since $a \mid a$ for all $a \in \mathbb{N}$. If $m \mid n$ and $n \mid m$, then $m = n$; hence, the relation is also antisymmetric. The relation is transitive, because if $m \mid n$ and $n \mid p$, then $m \mid p$.

Example 19.6. Let $X = \{1, 2, 3, 4, 6, 8, 12, 24\}$ be the set of divisors of 24 with the partial order defined in Example 19.5. Figure 19.7 shows the partial order on X.
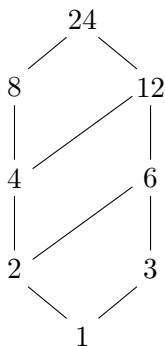
Figure 19.7: A partial order on the divisors of 24

Let Y be a subset of a poset X. An element u in X is an **_upper bound_** of Y if $a \preceq u$ for every element $a \in Y$. If u is an upper bound of Y such that $u \preceq v$ for every other upper bound v of Y, then u is called a **_least upper bound_** or **_supremum_** of Y. An element l in X is said to be a **_lower bound_** of Y if $l \preceq a$ for all $a \in Y$. If l is a lower bound of Y such that $k \preceq l$ for every other lower bound k of Y, then l is called a **_greatest lower bound_** or **_infimum_** of Y.

Example 19.8. Let $Y = \{2, 3, 4, 6\}$ be contained in the set X of Example 19.6. Then Y has upper bounds 12 and 24, with 12 as a least upper bound. The only lower bound is 1; hence, it must be a greatest lower bound.

As it turns out, least upper bounds and greatest lower bounds are unique if they exist.

Theorem 19.9. *Let Y be a nonempty subset of a poset X. If Y has a least upper bound, then Y has a unique least upper bound. If Y has a greatest lower bound, then Y has a unique greatest lower bound.*

PROOF.   Let u_1 and u_2 be least upper bounds for Y. By the definition of the least upper bound, $u_1 \preceq u$ for all upper bounds u of Y. In particular, $u_1 \preceq u_2$. Similarly, $u_2 \preceq u_1$. Therefore, $u_1 = u_2$ by antisymmetry. A similar argument show that the greatest lower bound is unique. □

On many posets it is possible to define binary operations by using the greatest lower bound and the least upper bound of two elements. A **lattice** is a poset L such that every pair of elements in L has a least upper bound and a greatest lower bound. The least upper bound of $a, b \in L$ is called the **join** of a and b and is denoted by $a \vee b$. The greatest lower bound of $a, b \in L$ is called the **meet** of a and b and is denoted by $a \wedge b$.

Example 19.10. Let X be a set. Then the power set of X, $\mathcal{P}(X)$, is a lattice. For two sets A and B in $\mathcal{P}(X)$, the least upper bound of A and B is $A \cup B$. Certainly $A \cup B$ is an upper bound of A and B, since $A \subset A \cup B$ and $B \subset A \cup B$. If C is some other set containing both A and B, then C must contain $A \cup B$; hence, $A \cup B$ is the least upper bound of A and B. Similarly, the greatest lower bound of A and B is $A \cap B$.

Example 19.11. Let G be a group and suppose that X is the set of subgroups of G. Then X is a poset ordered by set-theoretic inclusion, $\subset$. The set of subgroups of G is also a lattice. If H and K are subgroups of G, the greatest lower bound of H and K is $H \cap K$. The set $H \cup K$ may not be a subgroup of G. We leave it as an exercise to show that the least upper bound of H and K is the subgroup generated by $H \cup K$.

In set theory we have certain duality conditions. For example, by De Morgan's laws, any statement about sets that is true about $(A \cup B)'$ must also be true about $A' \cap B'$. We also have a duality principle for lattices.

Axiom 19.12 Principle of Duality. *Any statement that is true for all lattices remains true when $\preceq$ is replaced by $\succeq$ and $\vee$ and $\wedge$ are interchanged throughout the statement.*

The following theorem tells us that a lattice is an algebraic structure with two binary operations that satisfy certain axioms.

Theorem 19.13. *If L is a lattice, then the binary operations $\vee$ and $\wedge$ satisfy the following properties for $a, b, c \in L$.*

1. *Commutative laws: $a \vee b = b \vee a$ and $a \wedge b = b \wedge a$.*

2. *Associative laws: $a \vee (b \vee c) = (a \vee b) \vee c$ and $a \wedge (b \wedge c) = (a \wedge b) \wedge c$.*

3. *Idempotent laws: $a \vee a = a$ and $a \wedge a = a$.*

4. *Absorption laws: $a \vee (a \wedge b) = a$ and $a \wedge (a \vee b) = a$.*

PROOF.   By the Principle of Duality, we need only prove the first statement in each part.

(1) By definition $a \vee b$ is the least upper bound of $\{a, b\}$, and $b \vee a$ is the least upper bound of $\{b, a\}$; however, $\{a, b\} = \{b, a\}$.

(2) We will show that $a \vee (b \vee c)$ and $(a \vee b) \vee c$ are both least upper bounds of $\{a, b, c\}$. Let $d = a \vee b$. Then $c \preceq d \vee c = (a \vee b) \vee c$. We also know that

$$a \preceq a \vee b = d \preceq d \vee c = (a \vee b) \vee c.$$

A similar argument demonstrates that $b \preceq (a \vee b) \vee c$. Therefore, $(a \vee b) \vee c$ is an upper bound of $\{a, b, c\}$. We now need to show that $(a \vee b) \vee c$ is the least upper bound of $\{a, b, c\}$. Let u be some other upper bound of $\{a, b, c\}$. Then $a \preceq u$ and $b \preceq u$; hence, $d = a \vee b \preceq u$.

Since $c \preceq u$, it follows that $(a \vee b) \vee c = d \vee c \preceq u$. Therefore, $(a \vee b) \vee c$ must be the least upper bound of $\{a, b, c\}$. The argument that shows $a \vee (b \vee c)$ is the least upper bound of $\{a, b, c\}$ is the same. Consequently, $a \vee (b \vee c) = (a \vee b) \vee c$.

(3) The join of a and a is the least upper bound of $\{a\}$; hence, $a \vee a = a$.

(4) Let $d = a \wedge b$. Then $a \preceq a \vee d$. On the other hand, $d = a \wedge b \preceq a$, and so $a \vee d \preceq a$. Therefore, $a \vee (a \wedge b) = a$. $\qquad\square$

Given any arbitrary set L with operations $\vee$ and $\wedge$, satisfying the conditions of the previous theorem, it is natural to ask whether or not this set comes from some lattice. The following theorem says that this is always the case.

Theorem 19.14. *Let L be a nonempty set with two binary operations $\vee$ and $\wedge$ satisfying the commutative, associative, idempotent, and absorption laws. We can define a partial order on L by $a \preceq b$ if $a \vee b = b$. Furthermore, L is a lattice with respect to $\preceq$ if for all $a, b \in L$, we define the least upper bound and greatest lower bound of a and b by $a \vee b$ and $a \wedge b$, respectively.*

PROOF.   We first show that L is a poset under $\preceq$. Since $a \vee a = a$, $a \preceq a$ and $\preceq$ is reflexive. To show that $\preceq$ is antisymmetric, let $a \preceq b$ and $b \preceq a$. Then $a \vee b = b$ and $b \vee a = a$. By the commutative law, $b = a \vee b = b \vee a = a$. Finally, we must show that $\preceq$ is transitive. Let $a \preceq b$ and $b \preceq c$. Then $a \vee b = b$ and $b \vee c = c$. Thus,

$$a \vee c = a \vee (b \vee c) = (a \vee b) \vee c = b \vee c = c,$$

or $a \preceq c$.

To show that L is a lattice, we must prove that $a \vee b$ and $a \wedge b$ are, respectively, the least upper and greatest lower bounds of a and b. Since $a = (a \vee b) \wedge a = a \wedge (a \vee b)$, it follows that $a \preceq a \vee b$. Similarly, $b \preceq a \vee b$. Therefore, $a \vee b$ is an upper bound for a and b. Let u be any other upper bound of both a and b. Then $a \preceq u$ and $b \preceq u$. But $a \vee b \preceq u$ since

$$(a \vee b) \vee u = a \vee (b \vee u) = a \vee u = u.$$

The proof that $a \wedge b$ is the greatest lower bound of a and b is left as an exercise. $\qquad\square$

19.2   Boolean Algebras

Let us investigate the example of the power set, $\mathcal{P}(X)$, of a set X more closely. The power set is a lattice that is ordered by inclusion. By the definition of the power set, the largest element in $\mathcal{P}(X)$ is X itself and the smallest element is $\emptyset$, the empty set. For any set A in $\mathcal{P}(X)$, we know that $A \cap X = A$ and $A \cup \emptyset = A$. This suggests the following definition for lattices. An element I in a poset X is a ***largest element*** if $a \preceq I$ for all $a \in X$. An element O is a ***smallest element*** of X if $O \preceq a$ for all $a \in X$.

Let A be in $\mathcal{P}(X)$. Recall that the complement of A is

$$A' = X \setminus A = \{x : x \in X \text{ and } x \notin A\}.$$

We know that $A \cup A' = X$ and $A \cap A' = \emptyset$. We can generalize this example for lattices. A lattice L with a largest element I and a smallest element O is ***complemented*** if for each $a \in L$, there exists an a' such that $a \vee a' = I$ and $a \wedge a' = O$.

In a lattice L, the binary operations $\vee$ and $\wedge$ satisfy commutative and associative laws; however, they need not satisfy the distributive law

$$a \wedge (b \vee c) = (a \wedge b) \vee (a \wedge c);$$

however, in $\mathcal{P}(X)$ the distributive law is satisfied since

$$A \cap (B \cup C) = (A \cap B) \cup (A \cap C)$$

for $A, B, C \in \mathcal{P}(X)$. We will say that a lattice L is **distributive** if the following distributive law holds:

$$a \wedge (b \vee c) = (a \wedge b) \vee (a \wedge c)$$

for all $a, b, c \in L$.

Theorem 19.15. *A lattice L is distributive if and only if*

$$a \vee (b \wedge c) = (a \vee b) \wedge (a \vee c)$$

for all $a, b, c \in L$.

PROOF.   Let us assume that L is a distributive lattice.

$$
\begin{aligned}
a \vee (b \wedge c) &= [a \vee (a \wedge c)] \vee (b \wedge c) \\
&= a \vee [(a \wedge c) \vee (b \wedge c)] \\
&= a \vee [(c \wedge a) \vee (c \wedge b)] \\
&= a \vee [c \wedge (a \vee b)] \\
&= a \vee [(a \vee b) \wedge c] \\
&= [(a \vee b) \wedge a] \vee [(a \vee b) \wedge c] \\
&= (a \vee b) \wedge (a \vee c).
\end{aligned}
$$

The converse follows directly from the Duality Principle.   □

A **Boolean algebra** is a lattice B with a greatest element I and a smallest element O such that B is both distributive and complemented. The power set of X, $\mathcal{P}(X)$, is our prototype for a Boolean algebra. As it turns out, it is also one of the most important Boolean algebras. The following theorem allows us to characterize Boolean algebras in terms of the binary relations $\vee$ and $\wedge$ without mention of the fact that a Boolean algebra is a poset.

Theorem 19.16. *A set B is a Boolean algebra if and only if there exist binary operations $\vee$ and $\wedge$ on B satisfying the following axioms.*

1. *$a \vee b = b \vee a$ and $a \wedge b = b \wedge a$ for $a, b \in B$.*

2. *$a \vee (b \vee c) = (a \vee b) \vee c$ and $a \wedge (b \wedge c) = (a \wedge b) \wedge c$ for $a, b, c \in B$.*

3. *$a \wedge (b \vee c) = (a \wedge b) \vee (a \wedge c)$ and $a \vee (b \wedge c) = (a \vee b) \wedge (a \vee c)$ for $a, b, c \in B$.*

4. *There exist elements I and O such that $a \vee O = a$ and $a \wedge I = a$ for all $a \in B$.*

5. *For every $a \in B$ there exists an $a' \in B$ such that $a \vee a' = I$ and $a \wedge a' = O$.*

PROOF.   Let B be a set satisfying (1)–(5) in the theorem. One of the idempotent laws is satisfied since

$$
\begin{aligned}
a &= a \vee O \\
&= a \vee (a \wedge a') \\
&= (a \vee a) \wedge (a \vee a')
\end{aligned}
$$

$$= (a \vee a) \wedge I$$
$$= a \vee a.$$

Observe that

$$I \vee b = (b \vee b') \vee b = (b' \vee b) \vee b = b' \vee (b \vee b) = b' \vee b = I.$$

Consequently, the first of the two absorption laws holds, since

$$a \vee (a \wedge b) = (a \wedge I) \vee (a \wedge b)$$
$$= a \wedge (I \vee b)$$
$$= a \wedge I$$
$$= a.$$

The other idempotent and absorption laws are proven similarly. Since B also satisfies (1)–(3), the conditions of Theorem 19.14 are met; therefore, B must be a lattice. Condition (4) tells us that B is a distributive lattice.

For $a \in B$, $O \vee a = a$; hence, $O \preceq a$ and O is the smallest element in B. To show that I is the largest element in B, we will first show that $a \vee b = b$ is equivalent to $a \wedge b = a$. Since $a \vee I = a$ for all $a \in B$, using the absorption laws we can determine that

$$a \vee I = (a \wedge I) \vee I = I \vee (I \wedge a) = I$$

or $a \preceq I$ for all a in B. Finally, since we know that B is complemented by (5), B must be a Boolean algebra.

Conversely, suppose that B is a Boolean algebra. Let I and O be the greatest and least elements in B, respectively. If we define $a \vee b$ and $a \wedge b$ as least upper and greatest lower bounds of $\{a, b\}$, then B is a Boolean algebra by Theorem 19.14, Theorem 19.15, and our hypothesis. $\square$

Many other identities hold in Boolean algebras. Some of these identities are listed in the following theorem.

Theorem 19.17. *Let B be a Boolean algebra. Then*

1. *$a \vee I = I$ and $a \wedge O = O$ for all $a \in B$.*

2. *If $a \vee b = a \vee c$ and $a \wedge b = a \wedge c$ for $a, b, c \in B$, then $b = c$.*

3. *If $a \vee b = I$ and $a \wedge b = O$, then $b = a'$.*

4. *$(a')' = a$ for all $a \in B$.*

5. *$I' = O$ and $O' = I$.*

6. *$(a \vee b)' = a' \wedge b'$ and $(a \wedge b)' = a' \vee b'$ (De Morgan's Laws).*

PROOF. We will prove only (2). The rest of the identities are left as exercises. For $a \vee b = a \vee c$ and $a \wedge b = a \wedge c$, we have

$$b = b \vee (b \wedge a)$$
$$= b \vee (a \wedge b)$$
$$= b \vee (a \wedge c)$$
$$= (b \vee a) \wedge (b \vee c)$$

$$= (a \vee b) \wedge (b \vee c)$$
$$= (a \vee c) \wedge (b \vee c)$$
$$= (c \vee a) \wedge (c \vee b)$$
$$= c \vee (a \wedge b)$$
$$= c \vee (a \wedge c)$$
$$= c \vee (c \wedge a)$$
$$= c. \qquad \square$$

Finite Boolean Algebras

A Boolean algebra is a ***finite Boolean algebra*** if it contains a finite number of elements as a set. Finite Boolean algebras are particularly nice since we can classify them up to isomorphism.

Let B and C be Boolean algebras. A bijective map $\phi : B \to C$ is an ***isomorphism*** of Boolean algebras if

$$\phi(a \vee b) = \phi(a) \vee \phi(b)$$
$$\phi(a \wedge b) = \phi(a) \wedge \phi(b)$$

for all a and b in B.

We will show that any finite Boolean algebra is isomorphic to the Boolean algebra obtained by taking the power set of some finite set X. We will need a few lemmas and definitions before we prove this result. Let B be a finite Boolean algebra. An element $a \in B$ is an ***atom*** of B if $a \neq O$ and $a \wedge b = a$ for all nonzero $b \in B$. Equivalently, a is an atom of B if there is no nonzero $b \in B$ distinct from a such that $O \preceq b \preceq a$.

Lemma 19.18. *Let B be a finite Boolean algebra. If b is a nonzero element of B, then there is an atom a in B such that $a \preceq b$.*

PROOF.   If b is an atom, let $a = b$. Otherwise, choose an element b_1, not equal to O or b, such that $b_1 \preceq b$. We are guaranteed that this is possible since b is not an atom. If b_1 is an atom, then we are done. If not, choose b_2, not equal to O or b_1, such that $b_2 \preceq b_1$. Again, if b_2 is an atom, let $a = b_2$. Continuing this process, we can obtain a chain

$$O \preceq \cdots \preceq b_3 \preceq b_2 \preceq b_1 \preceq b.$$

Since B is a finite Boolean algebra, this chain must be finite. That is, for some k, b_k is an atom. Let $a = b_k$. $\qquad \square$

Lemma 19.19. *Let a and b be atoms in a finite Boolean algebra B such that $a \neq b$. Then $a \wedge b = O$.*

PROOF.   Since $a \wedge b$ is the greatest lower bound of a and b, we know that $a \wedge b \preceq a$. Hence, either $a \wedge b = a$ or $a \wedge b = O$. However, if $a \wedge b = a$, then either $a \preceq b$ or $a = O$. In either case we have a contradiction because a and b are both atoms; therefore, $a \wedge b = O$. $\qquad \square$

Lemma 19.20. *Let B be a Boolean algebra and $a, b \in B$. The following statements are equivalent.*

1. $a \preceq b$.

2. $a \wedge b' = O$.

3. $a' \vee b = I$.

PROOF. (1) $\Rightarrow$ (2). If $a \preceq b$, then $a \vee b = b$. Therefore,

$$
\begin{aligned}
a \wedge b' &= a \wedge (a \vee b)' \\
&= a \wedge (a' \wedge b') \\
&= (a \wedge a') \wedge b' \\
&= O \wedge b' \\
&= O.
\end{aligned}
$$

(2) $\Rightarrow$ (3). If $a \wedge b' = O$, then $a' \vee b = (a \wedge b')' = O' = I$.

(3) $\Rightarrow$ (1). If $a' \vee b = I$, then

$$
\begin{aligned}
a &= a \wedge (a' \vee b) \\
&= (a \wedge a') \vee (a \wedge b) \\
&= O \vee (a \wedge b) \\
&= a \wedge b.
\end{aligned}
$$

Thus, $a \preceq b$. $\qquad \square$

Lemma 19.21. *Let B be a Boolean algebra and b and c be elements in B such that $b \not\preceq c$. Then there exists an atom $a \in B$ such that $a \preceq b$ and $a \not\preceq c$.*

PROOF. By Lemma 19.20, $b \wedge c' \neq O$. Hence, there exists an atom a such that $a \preceq b \wedge c'$. Consequently, $a \preceq b$ and $a \not\preceq c$. $\qquad \square$

Lemma 19.22. *Let $b \in B$ and $a_1, \ldots, a_n$ be the atoms of B such that $a_i \preceq b$. Then $b = a_1 \vee \cdots \vee a_n$. Furthermore, if $a, a_1, \ldots, a_n$ are atoms of B such that $a \preceq b$, $a_i \preceq b$, and $b = a \vee a_1 \vee \cdots \vee a_n$, then $a = a_i$ for some $i = 1, \ldots, n$.*

PROOF. Let $b_1 = a_1 \vee \cdots \vee a_n$. Since $a_i \preceq b$ for each i, we know that $b_1 \preceq b$. If we can show that $b \preceq b_1$, then the lemma is true by antisymmetry. Assume $b \not\preceq b_1$. Then there exists an atom a such that $a \preceq b$ and $a \not\preceq b_1$. Since a is an atom and $a \preceq b$, we can deduce that $a = a_i$ for some a_i. However, this is impossible since $a \preceq b_1$. Therefore, $b \preceq b_1$.

Now suppose that $b = a_1 \vee \cdots \vee a_n$. If a is an atom less than b,

$$
a = a \wedge b = a \wedge (a_1 \vee \cdots \vee a_n) = (a \wedge a_1) \vee \cdots \vee (a \wedge a_n).
$$

But each term is O or a with $a \wedge a_i$ occurring for only one a_i. Hence, by Lemma 19.19, $a = a_i$ for some i. $\qquad \square$

Theorem 19.23. *Let B be a finite Boolean algebra. Then there exists a set X such that B is isomorphic to $\mathcal{P}(X)$.*

PROOF. We will show that B is isomorphic to $\mathcal{P}(X)$, where X is the set of atoms of B. Let $a \in B$. By Lemma 19.22, we can write a uniquely as $a = a_1 \vee \cdots \vee a_n$ for $a_1, \ldots, a_n \in X$. Consequently, we can define a map $\phi : B \to \mathcal{P}(X)$ by

$$
\phi(a) = \phi(a_1 \vee \cdots \vee a_n) = \{a_1, \ldots, a_n\}.
$$

Clearly, ϕ is onto.

Now let $a = a_1 \vee \cdots \vee a_n$ and $b = b_1 \vee \cdots \vee b_m$ be elements in B, where each a_i and each b_i is an atom. If $\phi(a) = \phi(b)$, then $\{a_1, \ldots, a_n\} = \{b_1, \ldots, b_m\}$ and $a = b$. Consequently, ϕ is injective.

The join of a and b is preserved by ϕ since

$$
\phi(a \vee b) = \phi(a_1 \vee \cdots \vee a_n \vee b_1 \vee \cdots \vee b_m)
$$

$$= \{a_1, \ldots, a_n, b_1, \ldots, b_m\}$$
$$= \{a_1, \ldots, a_n\} \cup \{b_1, \ldots, b_m\}$$
$$= \phi(a_1 \vee \cdots \vee a_n) \cup \phi(b_1 \wedge \cdots \vee b_m)$$
$$= \phi(a) \cup \phi(b).$$

Similarly, $\phi(a \wedge b) = \phi(a) \cap \phi(b)$. $\qquad\qquad\qquad\qquad\qquad\qquad\qquad\qquad\qquad\square$

We leave the proof of the following corollary as an exercise.

Corollary 19.24. *The order of any finite Boolean algebra must be 2^n for some positive integer n.*

19.3   The Algebra of Electrical Circuits

The usefulness of Boolean algebras has become increasingly apparent over the past several decades with the development of the modern computer. The circuit design of computer chips can be expressed in terms of Boolean algebras. In this section we will develop the Boolean algebra of electrical circuits and switches; however, these results can easily be generalized to the design of integrated computer circuitry.

A **switch** is a device, located at some point in an electrical circuit, that controls the flow of current through the circuit. Each switch has two possible states: it can be **open**, and not allow the passage of current through the circuit, or a it can be **closed**, and allow the passage of current. These states are mutually exclusive. We require that every switch be in one state or the other—a switch cannot be open and closed at the same time. Also, if one switch is always in the same state as another, we will denote both by the same letter; that is, two switches that are both labeled with the same letter a will always be open at the same time and closed at the same time.

Given two switches, we can construct two fundamental types of circuits. Two switches a and b are in **series** if they make up a circuit of the type that is illustrated in Figure 19.25. Current can pass between the terminals A and B in a series circuit only if both of the switches a and b are closed. We will denote this combination of switches by $a \wedge b$. Two switches a and b are in **parallel** if they form a circuit of the type that appears in Figure 19.26. In the case of a parallel circuit, current can pass between A and B if either one of the switches is closed. We denote a parallel combination of circuits a and b by $a \vee b$.

$$A \overline{\qquad} a \overline{\qquad} b \overline{\qquad} B$$
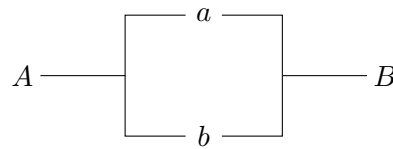
Figure 19.25: $a \wedge b$

Figure 19.26: $a \vee b$

We can build more complicated electrical circuits out of series and parallel circuits by replacing any switch in the circuit with one of these two fundamental types of circuits. Circuits constructed in this manner are called **series-parallel circuits**.

We will consider two circuits equivalent if they act the same. That is, if we set the switches in equivalent circuits exactly the same we will obtain the same result. For example,

in a series circuit $a \wedge b$ is exactly the same as $b \wedge a$. Notice that this is exactly the commutative law for Boolean algebras. In fact, the set of all series-parallel circuits forms a Boolean algebra under the operations of $\vee$ and $\wedge$. We can use diagrams to verify the different axioms of a Boolean algebra. The distributive law, $a \wedge (b \vee c) = (a \wedge b) \vee (a \wedge c)$, is illustrated in Figure 19.27. If a is a switch, then a' is the switch that is always open when a is closed and always closed when a is open. A circuit that is always closed is I in our algebra; a circuit that is always open is O. The laws for $a \wedge a' = O$ and $a \vee a' = I$ are shown in Figure 19.28.
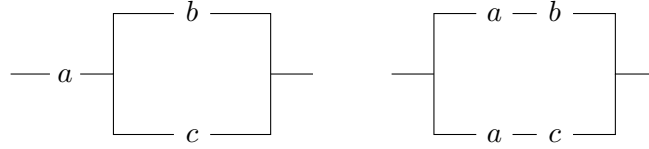
Figure 19.27: $a \wedge (b \vee c) = (a \wedge b) \vee (a \wedge c)$

Figure 19.28: $a \wedge a' = O$ and $a \vee a' = I$

Example 19.29. Every Boolean expression represents a switching circuit. For example, given the expression $(a \vee b) \wedge (a \vee b') \wedge (a \vee b)$, we can construct the circuit in Figure 19.32.

Theorem 19.30. *The set of all circuits is a Boolean algebra.*

We leave as an exercise the proof of this theorem for the Boolean algebra axioms not yet verified. We can now apply the techniques of Boolean algebras to switching theory.

Example 19.31. Given a complex circuit, we can now apply the techniques of Boolean algebra to reduce it to a simpler one. Consider the circuit in Figure 19.32. Since

$$
\begin{aligned}
(a \vee b) \wedge (a \vee b') \wedge (a \vee b) &= (a \vee b) \wedge (a \vee b) \wedge (a \vee b') \\
&= (a \vee b) \wedge (a \vee b') \\
&= a \vee (b \wedge b') \\
&= a \vee O \\
&= a,
\end{aligned}
$$

we can replace the more complicated circuit with a circuit containing the single switch a and achieve the same function.
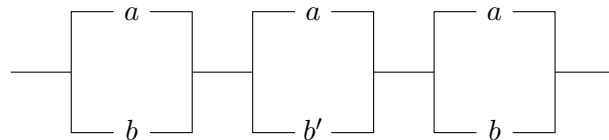
Figure 19.32: $(a \vee b) \wedge (a \vee b') \wedge (a \vee b)$

George Boole (1815–1864) was the first person to study lattices. In 1847, he published *The Investigation of the Laws of Thought*, a book in which he used lattices to formalize logic and the calculus of propositions. Boole believed that mathematics was the study of form rather than of content; that is, he was not so much concerned with what he was calculating as with how he was calculating it. Boole's work was carried on by his friend Augustus De Morgan (1806–1871). De Morgan observed that the principle of duality often held in set theory, as is illustrated by De Morgan's laws for set theory. He believed, as did Boole, that mathematics was the study of symbols and abstract operations.

Set theory and logic were further advanced by such mathematicians as Alfred North Whitehead (1861–1947), Bertrand Russell (1872–1970), and David Hilbert (1862–1943). In *Principia Mathematica*, Whitehead and Russell attempted to show the connection between mathematics and logic by the deduction of the natural number system from the rules of formal logic. If the natural numbers could be determined from logic itself, then so could much of the rest of existing mathematics. Hilbert attempted to build up mathematics by using symbolic logic in a way that would prove the consistency of mathematics. His approach was dealt a mortal blow by Kurt Gödel (1906–1978), who proved that there will always be "undecidable" problems in any sufficiently rich axiomatic system; that is, that in any mathematical system of any consequence, there will always be statements that can never be proven either true or false.

As often occurs, this basic research in pure mathematics later became indispensable in a wide variety of applications. Boolean algebras and logic have become essential in the design of the large-scale integrated circuitry found on today's computer chips. Sociologists have used lattices and Boolean algebras to model social hierarchies; biologists have used them to describe biosystems.

19.4 Exercises

1.   Draw the lattice diagram for the power set of $X = \{a, b, c, d\}$ with the set inclusion relation, $\subset$.

2.   Draw the diagram for the set of positive integers that are divisors of 30. Is this poset a Boolean algebra?

3.   Draw a diagram of the lattice of subgroups of $\mathbb{Z}_{12}$.

4.   Let B be the set of positive integers that are divisors of 36. Define an order on B by $a \preceq b$ if $a \mid b$. Prove that B is a Boolean algebra. Find a set X such that B is isomorphic to $\mathcal{P}(X)$.

5.   Prove or disprove: $\mathbb{Z}$ is a poset under the relation $a \preceq b$ if $a \mid b$.

6.   Draw the switching circuit for each of the following Boolean expressions.

(a) $(a \vee b \vee a') \wedge a$ (c) $a \vee (a \wedge b)$

(b) $(a \vee b)' \wedge (a \vee b)$ (d) $(c \vee a \vee b) \wedge c' \wedge (a \vee b)'$

7.   Draw a circuit that will be closed exactly when only one of three switches a, b, and c are closed.

8.   Prove or disprove that the two circuits shown are equivalent.
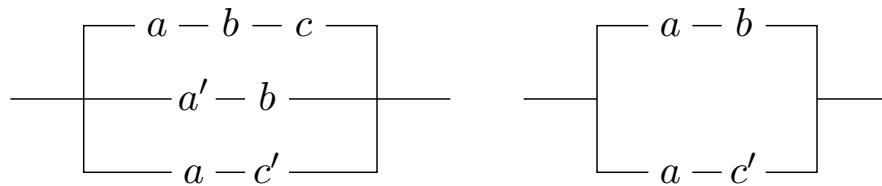
9. Let X be a finite set containing n elements. Prove that $\mathcal{P}(X) = 2^n$. Conclude that the order of any finite Boolean algebra must be 2^n for some $n \in \mathbb{N}$.

10. For each of the following circuits, write a Boolean expression. If the circuit can be replaced by one with fewer switches, give the Boolean expression and draw a diagram for the new circuit.
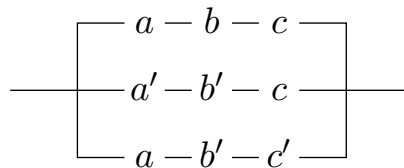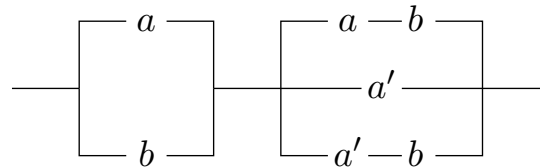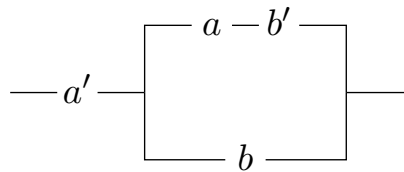
11. Prove or disprove: The set of all nonzero integers is a lattice, where $a \preceq b$ is defined by $a \mid b$.

12. Let L be a nonempty set with two binary operations $\vee$ and $\wedge$ satisfying the commutative, associative, idempotent, and absorption laws. We can define a partial order on L, as in Theorem 19.14, by $a \preceq b$ if $a \vee b = b$. Prove that the greatest lower bound of a and b is $a \wedge b$.

13. Let G be a group and X be the set of subgroups of G ordered by set-theoretic inclusion. If H and K are subgroups of G, show that the least upper bound of H and K is the subgroup generated by $H \cup K$.

14. Let R be a ring and suppose that X is the set of ideals of R. Show that X is a poset ordered by set-theoretic inclusion, $\subset$. Define the meet of two ideals I and J in X by $I \cap J$ and the join of I and J by $I + J$. Prove that the set of ideals of R is a lattice under these operations.

15. Let B be a Boolean algebra. Prove each of the following identities.

(a) $a \vee I = I$ and $a \wedge O = O$ for all $a \in B$.

(b) If $a \vee b = I$ and $a \wedge b = O$, then $b = a'$.

(c) $(a')' = a$ for all $a \in B$.

(d) $I' = O$ and $O' = I$.

(e) $(a \vee b)' = a' \wedge b'$ and $(a \wedge b)' = a' \vee b'$ (De Morgan's laws).

16.   By drawing the appropriate diagrams, complete the proof of Theorem 19.30 to show that the switching functions form a Boolean algebra.

17.   Let B be a Boolean algebra. Define binary operations $+$ and $\cdot$ on B by

$$a + b = (a \wedge b') \vee (a' \wedge b)$$
$$a \cdot b = a \wedge b.$$

Prove that B is a commutative ring under these operations satisfying $a^2 = a$ for all $a \in B$.

18.   Let X be a poset such that for every a and b in X, either $a \preceq b$ or $b \preceq a$. Then X is said to be a **totally ordered set**.

(a) Is $a \mid b$ a total order on $\mathbb{N}$?

(b) Prove that $\mathbb{N}$, $\mathbb{Z}$, $\mathbb{Q}$, and $\mathbb{R}$ are totally ordered sets under the usual ordering $\leq$.

19.   Let X and Y be posets. A map $\phi : X \to Y$ is **order-preserving** if $a \preceq b$ implies that $\phi(a) \preceq \phi(b)$. Let L and M be lattices. A map $\psi : L \to M$ is a **lattice homomorphism** if $\psi(a \vee b) = \psi(a) \vee \psi(b)$ and $\psi(a \wedge b) = \psi(a) \wedge \psi(b)$. Show that every lattice homomorphism is order-preserving, but that it is not the case that every order-preserving homomorphism is a lattice homomorphism.

20.   Let B be a Boolean algebra. Prove that $a = b$ if and only if $(a \wedge b') \vee (a' \wedge b) = O$ for $a, b \in B$.

21.   Let B be a Boolean algebra. Prove that $a = O$ if and only if $(a \wedge b') \vee (a' \wedge b) = b$ for all $b \in B$.

22.   Let L and M be lattices. Define an order relation on $L \times M$ by $(a, b) \preceq (c, d)$ if $a \preceq c$ and $b \preceq d$. Show that $L \times M$ is a lattice under this partial order.

19.5 Programming Exercises

1.   A **Boolean** or **switching function on** n **variables** is a map $f : \{O, I\}^n \to \{0, I\}$. A Boolean polynomial is a special type of Boolean function: it is any type of Boolean expression formed from a finite combination of variables $x_1, \ldots, x_n$ together with O and I, using the operations $\vee$, $\wedge$, and $'$. The values of the functions are defined in Table 19.33. Write a program to evaluate Boolean polynomials.

x	y	x'	$x \vee y$	$x \wedge y$
0	0	1	0	0
0	1	1	1	0
1	0	0	1	0
1	1	0	1	1

Table 19.33: Boolean polynomials

References and Suggested Readings

[1] Donnellan, T. *Lattice Theory.* Pergamon Press, Oxford, 1968.

[2] Halmos, P. R. "The Basic Concepts of Algebraic Logic," *American Mathematical Monthly* **53**(1956), 363–87.

[3] Hohn, F. "Some Mathematical Aspects of Switching," *American Mathematical Monthly* **62**(1955), 75–90.

[4] Hohn, F. *Applied Boolean Algebra.* 2nd ed. Macmillan, New York, 1966.

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[6] Whitesitt, J. *Boolean Algebra and Its Applications.* Dover, Mineola, NY, 2010.

19.6 Sage

Sage has support for both partially ordered sets ("posets") and lattices, and does an excellent job of providing visual depictions of both.

Creating Partially Ordered Sets

Example 19.6 in the text is a good example to replicate as a demonstration of Sage commands. We first define the elements of the set X.

```
X = (24).divisors()
X
```

```
[1, 2, 3, 4, 6, 8, 12, 24]
```

One approach to creating the relation is to specify *every* instance where one element is comparable to the another. So we build a list of pairs, where each pair contains comparable elements, with the lesser one first. This is the set of relations.

```
R = [(a,b) for a in X for b in X if a.divides(b)]; R
```

```
[(1, 1), (1, 2), (1, 3), (1, 4), (1, 6), (1, 8), (1, 12), (1, 24),
 (2, 2), (2, 4), (2, 6), (2, 8), (2, 12), (2, 24), (3, 3), (3, 6),
 (3, 12), (3, 24), (4, 4), (4, 8), (4, 12), (4, 24), (6, 6),
 (6, 12), (6, 24), (8, 8), (8, 24), (12, 12), (12, 24), (24, 24)]
```

We construct the poset by giving the the Poset constructor a list containing the elements and the relations. We can then easily get a "plot" of the poset. Notice the plot just shows the "cover relations" — a minimal set of comparisons which the assumption of transitivity would expand into the set of all the relations.

```
D = Poset([X, R])
D.plot()
```

Another approach to creating a Poset is to let the poset constructor run over all the pairs of elements, and all we do is give the constructor a way to test if two elements are comparable. Our comparison function should expect two elements and then return True or False. A "lambda" function is one way to quickly build such a function. This may be a new idea for you, but mastering lambda functions can be a great convenience. Notice that "lambda" is a

word reserved for just this purpose (so, for example, `lambda` is a bad choice for the name of
an eigenvalue of a matrix). There are other ways to make functions in Sage, but a lambda
function is quickest when the function is simple.

```
divisible = lambda x, y: x.divides(y)
L = Poset([X, divisible])
L == D
```

True

```
L.plot()
```

Sage also has a collection of stock posets. Some are one-shot constructions, while others
are members of parameterized families. Use tab-completion on `Posets.` to see the full list.
Here are some examples.

A one-shot construction. Perhaps what you would expect, though there might be other,
equally plausible, alternatives.

```
Q = Posets.PentagonPoset()
Q.plot()
```

A parameterized family. This is the classic example where the elements are subsets of a set
with n elements and the relation is "subset of."

```
S = Posets.BooleanLattice(4)
S.plot()
```

And random posets. These can be useful for testing and experimenting, but are unlikely to
exhibit special cases that may be important. You might run the following command many
times and vary the second argument, which is a rough upper bound on the probability any
two elements are comparable. Remember that the plot only shows the cover relations. The
more elements that are comparable, the more "vertically stretched" the plot will be.

```
T = Posets.RandomPoset(20,0.05)
T.plot()
```

Properties of a Poset

Once you have a poset, what can you do with it? Let's return to our first example, `D`.
We can of course determine if one element is less than another, which is the fundamental
structure of a poset.

```
D.is_lequal(4, 8)
```

True

```
D.is_lequal(4, 4)
```

True

```
D.is_less_than(4, 8)
```

True

```
D.is_less_than(4, 4)
```

```
False
```

```
D.is_lequal(6, 8)
```

```
False
```

```
D.is_lequal(8, 6)
```

```
False
```

Notice that 6 and 8 are not comparable in this poset (it is a *partial* order). The methods `.is_gequal()` and `.is_greater_than()` work similarly, but returns `True` if the first element is greater (or equal).

```
D.is_gequal(8, 4)
```

```
True
```

```
D.is_greater_than(4, 8)
```

```
False
```

We can find the largest and smallest elements of a poset. This is a random poset built with a 10%probability, but copied here to be repeatable.

```
X = range(20)
C = [[18, 7],   [9, 11], [9, 10], [11, 8], [6, 10],
     [10, 2],   [0, 2],   [2, 1],   [1, 8], [8, 12],
     [8, 3],   [3, 15], [15, 7], [7, 16],   [7, 4],
     [16, 17], [16, 13], [4, 19], [4, 14], [14, 5]]
P = Poset([X, C])
P.plot()
```

```
P.minimal_elements()
```

```
[18, 9, 6, 0]
```

```
P.maximal_elements()
```

```
[5, 19, 13, 17, 12]
```

Elements of a poset can be partioned into level sets. In plots of posets, elements at the same level are plotted vertically at the same height. Each level set is obtained by removing all of the previous level sets and then taking the minimal elements of the result.

```
P.level_sets()
```

```
[[18, 9, 6, 0], [11, 10], [2], [1], [8], [3, 12],
 [15], [7], [16, 4], [13, 17, 14, 19], [5]]
```

If we make two elements in `R` comparable when they had not previously been, this is an extension of `R`. Consider all possible extensions of one poset — we can make a poset from all of these, where set inclusion is the relation. A linear extension is a maximal element in this poset of posets. Informally, we are adding as many new relations as possible, consistent with the original poset and so that the result is a total order. In other words, there is an ordering of the elements that is consistent with the order in the poset. We can build such a thing, but the output is just a list of the elements in the linear order. A computer scientist would be inclined to call this a "topological sort."

```
linear = P.linear_extension(); linear
```

```
[18, 9, 11, 6, 10, 0, 2, 1, 8, 3, 15,
 7, 4, 14, 5, 19, 16, 13, 17, 12]
```

We can construct subposets by giving a set of elements to induce the new poset. Here we take roughly the "bottom half" of the random poset P by inducing the subposet on a union of some of the level sets.

```
level = P.level_sets()
bottomhalf = sum([level[i] for i in range(5)], [])
B = P.subposet(bottomhalf)
B.plot()
```

The dual of a poset retains the same set of elements, but reverses any comparisons.

```
Pdual = P.dual()
Pdual.plot()
```

Taking the dual of the divisibility poset from Example 19.6 would be like changing the relation to "is a multiple of."

```
Ddual = D.dual()
Ddual.plot()
```

Lattices

Every lattice is a poset, so all the commands above will perform equally well for a lattice. But how do you create a lattice? Simple — first create a poset and then feed it into the `LatticePoset()` constructor. But realize that just because you give this constructor a poset, it does not mean a lattice will always come back out. Only if the poset is *already* a lattice will it get upgraded from a poset to a lattice for Sage's purposes, and you will get a `ValueError` if the upgrade is not possible. Finally, notice that some of the posets Sage constructs are already recognized as lattices, such as the prototypical `BooleanLattice`.

```
P = Posets.AntichainPoset(8)
P.is_lattice()
```

```
False
```

```
LatticePoset(P)
```

```
Traceback (most recent call last):
...
ValueError: not a meet-semilattice: no bottom element
```

An integer composition of n is a list of positive integers that sum to n. A composition C_1 covers a composition C_2 if C_2 can be formed from C_1 by adding consecutive parts. For example, $C_1 = [2, 1, 2] \succeq [3, 2] = C_2$. With this relation, the set of all integer compositions of a fixed integer n is a poset that is also a lattice.

```
CP = Posets.IntegerCompositions(5)
C = LatticePoset(CP)
C.plot()
```

A meet or a join is a fundamental operation in a lattice.

```
par = C.an_element().parent()
a = par([1, 1, 1, 2])
b = par([2, 1, 1, 1])
a, b
```

([1, 1, 1, 2], [2, 1, 1, 1])

```
C.meet(a, b)
```

[2, 1, 2]

```
c = par([1, 4])
d = par([2, 3])
c, d
```

([1, 4], [2, 3])

```
C.join(c, d)
```

[1, 1, 3]

Once a poset is upgraded to lattice status, then additional commands become available, or the character of their results changes.

An example of the former is the `.is_distributive()` method.

```
C.is_distributive()
```

True

An example of the latter is the `.top()` method. What your text calls a largest element and a smallest element of a lattice, Sage calls a top and a bottom. For a poset, `.top()` and `.bottom()` may return an element or may not (returning `None`), but for a lattice it is guaranteed to return exactly one element.

```
C.top()
```

[1, 1, 1, 1, 1]

```
C.bottom()
```

[5]

Notice that the returned values are all elements of the lattice, in this case ordered lists of integers summing to 5.

Complements now make sense in a lattice. The result of the `.complements()` method is a dictionary that uses elements of the lattice as the keys. We say the dictionary is "indexed" by the elements of the lattice. The result is a list of the complements of the element. We call this the "value" of the key-value pair. (You may know dictionaries as "associative arrays", but they are really just fancy functions.)

```
comp = C.complements()
comp[par([1, 1, 1, 2])]
```

[[4, 1]]

The lattice of integer compositions is a complemented lattice, as we can see by the result that each element has a single (unique) complement, evidenced by the lists of length 1 in the values of the dictionary. Or we can just ask Sage via `.is_complemented()`. Dictionaries have no inherent order, so you may get different output each time you inspect the dictionary.

```
comp
```

```
{[1, 1, 1, 1, 1]: [[5]],
 [1, 1, 1, 2]: [[4, 1]],
 [1, 1, 2, 1]: [[3, 2]],
 [1, 1, 3]: [[3, 1, 1]],
 [1, 2, 1, 1]: [[2, 3]],
 [1, 2, 2]: [[2, 2, 1]],
 [1, 3, 1]: [[2, 1, 2]],
 [1, 4]: [[2, 1, 1, 1]],
 [2, 1, 1, 1]: [[1, 4]],
 [2, 1, 2]: [[1, 3, 1]],
 [2, 2, 1]: [[1, 2, 2]],
 [2, 3]: [[1, 2, 1, 1]],
 [3, 1, 1]: [[1, 1, 3]],
 [3, 2]: [[1, 1, 2, 1]],
 [4, 1]: [[1, 1, 1, 2]],
 [5]: [[1, 1, 1, 1, 1]]}
```

```
[len(e[1]) for e in comp.items()]
```

```
[1, 1, 1, 1, 1, 1, 1, 1, 1, 1, 1, 1, 1, 1, 1, 1]
```

```
C.is_complemented()
```

```
True
```

There are many more commands which apply to posets and lattices, so build a few and use tab-completion liberally to explore. There is more to discover than we can cover in just a single chapter, but you now have the basic tools to profitably study posets and lattices in Sage.

19.7   Sage Exercises

1.   Use `R = Posets.RandomPoset(30,0.05)` to construct a random poset. Use `R.plot()` to get an idea of what you have built.

  (a) Illustrate the use of the poset methods: `.is_lequal()`, `.is_less_than()`, `.is_gequal()`, and `.is_greater_than()` to determine if two specific elements (of your choice) are related or incomparable.

  (b) Use `.minimal_elements()` and `.maximal_elements()` to find the smallest and largest elements of your poset.

  (c) Use `LatticePoset(R)` to see if the poset `R` is a lattice by attempting to convert it into a lattice.

  (d) Find a linear extension of your poset. Confirm that any pair of elements that are comparable in the poset will be similarly comparable in the linear extension.

2.   Construct the poset on the positive divisors of $72 = 2^3 \cdot 3^2$ with divisiblity as the relation, and then convert to a lattice.

  (a) Determine the one and zero element using `.top()` and `.bottom()`.

(b) Determine all the pairs of elements of the lattice that are complements of each other *without* using the `.complement()` method, but rather just use the `.meet()` and `.join()` methods. Extra credit if you can output each pair just once.

(c) Determine if the lattice is distributive using just the `.meet()` and `.join()` methods, and not the `.is_distributive()` method.

3. Construct several specific diamond lattices with `Posets.DiamondPoset(n)` by varying the value of `n`. Once you feel you have enough empirical evidence, give answers, with justifications, to the following questions for *general* values of n, based on observations obtained from your experiments with Sage.

(a) Which elements have complements and which do not, and why?

(b) Read the documentation of the `.antichains()` method to learn what an antichain is. How many antichains are there?

(c) Is the lattice distributive?

4. Use `Posets.BooleanLattice(4)` to construct an instance of the prototypical Boolean algebra on 16 elements (i.e., all subsets of a 4-set).

Then use `Posets.IntegerCompositions(5)` to construct the poset whose 16 elements are the compositions of the integer 5. We have seen above that the integer composition lattice is distributive and complemented, making it a Boolean algebra. And by Theorem 19.23 we can conclude that these two Boolean algebras are isomorphic.

Use the `.plot()` method to see the similarity visually. Then use the method `.hasse_diagram()` on each poset to obtain a directed graph (which you can also plot, though the embedding into the plane may not be as informative). Employ the graph method `.is_isomorphic()` to see that the two Hasse diagrams really are the "same."

5. (Advanced) For the previous question, construct an *explicit* isomorphism between the two Boolean algebras. This would be a bijective function (constructed with the `def` command) that converts compositions into sets (or if, you choose, sets into compositions) and which respects the meet and join operations. You can test and illustrate your function by its interaction with specific elements evaluated in the meet and join operations, as described in the definition of an isomorphism of Boolean algebras.

$$\boxed{20}$$

Vector Spaces

In a physical system a quantity can often be described with a single number. For example, we need to know only a single number to describe temperature, mass, or volume. However, for some quantities, such as location, we need several numbers. To give the location of a point in space, we need x, y, and z coordinates. Temperature distribution over a solid object requires four numbers: three to identify each point within the object and a fourth to describe the temperature at that point. Often n-tuples of numbers, or vectors, also have certain algebraic properties, such as addition or scalar multiplication.

In this chapter we will examine mathematical structures called vector spaces. As with groups and rings, it is desirable to give a simple list of axioms that must be satisfied to make a set of vectors a structure worth studying.

20.1   Definitions and Examples

A ***vector space*** V over a field F is an abelian group with a ***scalar product*** $\alpha \cdot v$ or αv defined for all $\alpha \in F$ and all $v \in V$ satisfying the following axioms.

- $\alpha(\beta v) = (\alpha\beta)v$;

- $(\alpha + \beta)v = \alpha v + \beta v$;

- $\alpha(u + v) = \alpha u + \alpha v$;

- $1v = v$;

where $\alpha, \beta \in F$ and $u, v \in V$.

The elements of V are called ***vectors***; the elements of F are called ***scalars***. It is important to notice that in most cases two vectors cannot be multiplied. In general, it is only possible to multiply a vector with a scalar. To differentiate between the scalar zero and the vector zero, we will write them as 0 and $\mathbf{0}$, respectively.

Let us examine several examples of vector spaces. Some of them will be quite familiar; others will seem less so.

Example 20.1. The n-tuples of real numbers, denoted by $\mathbb{R}^n$, form a vector space over $\mathbb{R}$. Given vectors $u = (u_1, \ldots, u_n)$ and $v = (v_1, \ldots, v_n)$ in $\mathbb{R}^n$ and α in $\mathbb{R}$, we can define vector addition by

$$u + v = (u_1, \ldots, u_n) + (v_1, \ldots, v_n) = (u_1 + v_1, \ldots, u_n + v_n)$$

and scalar multiplication by

$$\alpha u = \alpha(u_1, \ldots, u_n) = (\alpha u_1, \ldots, \alpha u_n).$$

Example 20.2. If F is a field, then $F[x]$ is a vector space over F. The vectors in $F[x]$ are simply polynomials, and vector addition is just polynomial addition. If $\alpha \in F$ and $p(x) \in F[x]$, then scalar multiplication is defined by $\alpha p(x)$.

Example 20.3. The set of all continuous real-valued functions on a closed interval $[a, b]$ is a vector space over $\mathbb{R}$. If $f(x)$ and $g(x)$ are continuous on $[a, b]$, then $(f + g)(x)$ is defined to be $f(x) + g(x)$. Scalar multiplication is defined by $(\alpha f)(x) = \alpha f(x)$ for $\alpha \in \mathbb{R}$. For example, if $f(x) = \sin x$ and $g(x) = x^2$, then $(2f + 5g)(x) = 2\sin x + 5x^2$.

Example 20.4. Let $V = \mathbb{Q}(\sqrt{2}) = \{a + b\sqrt{2} : a, b \in \mathbb{Q}\}$. Then V is a vector space over $\mathbb{Q}$. If $u = a + b\sqrt{2}$ and $v = c + d\sqrt{2}$, then $u + v = (a + c) + (b + d)\sqrt{2}$ is again in V. Also, for $\alpha \in \mathbb{Q}$, αv is in V. We will leave it as an exercise to verify that all of the vector space axioms hold for V.

Proposition 20.5. *Let V be a vector space over F. Then each of the following statements is true.*

1. $0v = \mathbf{0}$ *for all $v \in V$.*

2. $\alpha\mathbf{0} = \mathbf{0}$ *for all $\alpha \in F$.*

3. *If $\alpha v = \mathbf{0}$, then either $\alpha = 0$ or $v = \mathbf{0}$.*

4. $(-1)v = -v$ *for all $v \in V$.*

5. $-(\alpha v) = (-\alpha)v = \alpha(-v)$ *for all $\alpha \in F$ and all $v \in V$.*

PROOF. To prove (1), observe that

$$0v = (0 + 0)v = 0v + 0v;$$

consequently, $\mathbf{0} + 0v = 0v + 0v$. Since V is an abelian group, $\mathbf{0} = 0v$.

The proof of (2) is almost identical to the proof of (1). For (3), we are done if $\alpha = 0$. Suppose that $\alpha \neq 0$. Multiplying both sides of $\alpha v = \mathbf{0}$ by $1/\alpha$, we have $v = \mathbf{0}$.

To show (4), observe that

$$v + (-1)v = 1v + (-1)v = (1 - 1)v = 0v = \mathbf{0},$$

and so $-v = (-1)v$. We will leave the proof of (5) as an exercise. $\square$

20.2 Subspaces

Just as groups have subgroups and rings have subrings, vector spaces also have substructures. Let V be a vector space over a field F, and W a subset of V. Then W is a **subspace** of V if it is closed under vector addition and scalar multiplication; that is, if $u, v \in W$ and $\alpha \in F$, it will always be the case that $u + v$ and αv are also in W.

Example 20.6. Let W be the subspace of $\mathbb{R}^3$ defined by $W = \{(x_1, 2x_1 + x_2, x_1 - x_2) : x_1, x_2 \in \mathbb{R}\}$. We claim that W is a subspace of $\mathbb{R}^3$. Since

$$\alpha(x_1, 2x_1 + x_2, x_1 - x_2) = (\alpha x_1, \alpha(2x_1 + x_2), \alpha(x_1 - x_2))$$
$$= (\alpha x_1, 2(\alpha x_1) + \alpha x_2, \alpha x_1 - \alpha x_2),$$

W is closed under scalar multiplication. To show that W is closed under vector addition, let $u = (x_1, 2x_1 + x_2, x_1 - x_2)$ and $v = (y_1, 2y_1 + y_2, y_1 - y_2)$ be vectors in W. Then

$$u + v = (x_1 + y_1, 2(x_1 + y_1) + (x_2 + y_2), (x_1 + y_1) - (x_2 + y_2)).$$

Example 20.7. Let W be the subset of polynomials of $F[x]$ with no odd-power terms. If $p(x)$ and $q(x)$ have no odd-power terms, then neither will $p(x) + q(x)$. Also, $\alpha p(x) \in W$ for $\alpha \in F$ and $p(x) \in W$.

Let V be any vector space over a field F and suppose that $v_1, v_2, \ldots, v_n$ are vectors in V and $\alpha_1, \alpha_2, \ldots, \alpha_n$ are scalars in F. Any vector w in V of the form

$$w = \sum_{i=1}^{n} \alpha_i v_i = \alpha_1 v_1 + \alpha_2 v_2 + \cdots + \alpha_n v_n$$

is called a *linear combination* of the vectors $v_1, v_2, \ldots, v_n$. The *spanning set* of vectors $v_1, v_2, \ldots, v_n$ is the set of vectors obtained from all possible linear combinations of $v_1, v_2, \ldots, v_n$. If W is the spanning set of $v_1, v_2, \ldots, v_n$, then we say that W is *spanned* by $v_1, v_2, \ldots, v_n$.

Proposition 20.8. *Let $S = \{v_1, v_2, \ldots, v_n\}$ be vectors in a vector space V. Then the span of S is a subspace of V.*

PROOF.   Let u and v be in S. We can write both of these vectors as linear combinations of the v_i's:

$$u = \alpha_1 v_1 + \alpha_2 v_2 + \cdots + \alpha_n v_n$$
$$v = \beta_1 v_1 + \beta_2 v_2 + \cdots + \beta_n v_n.$$

Then

$$u + v = (\alpha_1 + \beta_1)v_1 + (\alpha_2 + \beta_2)v_2 + \cdots + (\alpha_n + \beta_n)v_n$$

is a linear combination of the v_i's. For $\alpha \in F$,

$$\alpha u = (\alpha\alpha_1)v_1 + (\alpha\alpha_2)v_2 + \cdots + (\alpha\alpha_n)v_n$$

is in the span of S. $\square$

20.3   Linear Independence

Let $S = \{v_1, v_2, \ldots, v_n\}$ be a set of vectors in a vector space V. If there exist scalars $\alpha_1, \alpha_2 \ldots \alpha_n \in F$ such that not all of the α_i's are zero and

$$\alpha_1 v_1 + \alpha_2 v_2 + \cdots + \alpha_n v_n = \mathbf{0},$$

then S is said to be *linearly dependent*. If the set S is not linearly dependent, then it is said to be *linearly independent*. More specifically, S is a linearly independent set if

$$\alpha_1 v_1 + \alpha_2 v_2 + \cdots + \alpha_n v_n = \mathbf{0}$$

implies that

$$\alpha_1 = \alpha_2 = \cdots = \alpha_n = 0$$

for any set of scalars $\{\alpha_1, \alpha_2 \ldots \alpha_n\}$.

Proposition 20.9. *Let $\{v_1, v_2, \ldots, v_n\}$ be a set of linearly independent vectors in a vector space. Suppose that*

$$v = \alpha_1 v_1 + \alpha_2 v_2 + \cdots + \alpha_n v_n = \beta_1 v_1 + \beta_2 v_2 + \cdots + \beta_n v_n.$$

Then $\alpha_1 = \beta_1, \alpha_2 = \beta_2, \ldots, \alpha_n = \beta_n$.

PROOF. If
$$v = \alpha_1 v_1 + \alpha_2 v_2 + \cdots + \alpha_n v_n = \beta_1 v_1 + \beta_2 v_2 + \cdots + \beta_n v_n,$$
then
$$(\alpha_1 - \beta_1)v_1 + (\alpha_2 - \beta_2)v_2 + \cdots + (\alpha_n - \beta_n)v_n = \mathbf{0}.$$
Since $v_1, \ldots, v_n$ are linearly independent, $\alpha_i - \beta_i = 0$ for $i = 1, \ldots, n$. □

The definition of linear dependence makes more sense if we consider the following proposition.

Proposition 20.10. *A set $\{v_1, v_2, \ldots, v_n\}$ of vectors in a vector space V is linearly dependent if and only if one of the v_i's is a linear combination of the rest.*

PROOF. Suppose that $\{v_1, v_2, \ldots, v_n\}$ is a set of linearly dependent vectors. Then there exist scalars $\alpha_1, \ldots, \alpha_n$ such that
$$\alpha_1 v_1 + \alpha_2 v_2 + \cdots + \alpha_n v_n = \mathbf{0},$$
with at least one of the α_i's not equal to zero. Suppose that $\alpha_k \neq 0$. Then
$$v_k = -\frac{\alpha_1}{\alpha_k}v_1 - \cdots - \frac{\alpha_{k-1}}{\alpha_k}v_{k-1} - \frac{\alpha_{k+1}}{\alpha_k}v_{k+1} - \cdots - \frac{\alpha_n}{\alpha_k}v_n.$$

Conversely, suppose that
$$v_k = \beta_1 v_1 + \cdots + \beta_{k-1}v_{k-1} + \beta_{k+1}v_{k+1} + \cdots + \beta_n v_n.$$
Then
$$\beta_1 v_1 + \cdots + \beta_{k-1}v_{k-1} - v_k + \beta_{k+1}v_{k+1} + \cdots + \beta_n v_n = \mathbf{0}. \quad □$$

The following proposition is a consequence of the fact that any system of homogeneous linear equations with more unknowns than equations will have a nontrivial solution. We leave the details of the proof for the end-of-chapter exercises.

Proposition 20.11. *Suppose that a vector space V is spanned by n vectors. If $m > n$, then any set of m vectors in V must be linearly dependent.*

A set $\{e_1, e_2, \ldots, e_n\}$ of vectors in a vector space V is called a **basis** for V if $\{e_1, e_2, \ldots, e_n\}$ is a linearly independent set that spans V.

Example 20.12. The vectors $e_1 = (1, 0, 0)$, $e_2 = (0, 1, 0)$, and $e_3 = (0, 0, 1)$ form a basis for $\mathbb{R}^3$. The set certainly spans $\mathbb{R}^3$, since any arbitrary vector (x_1, x_2, x_3) in $\mathbb{R}^3$ can be written as $x_1 e_1 + x_2 e_2 + x_3 e_3$. Also, none of the vectors e_1, e_2, e_3 can be written as a linear combination of the other two; hence, they are linearly independent. The vectors e_1, e_2, e_3 are not the only basis of $\mathbb{R}^3$: the set $\{(3, 2, 1), (3, 2, 0), (1, 1, 1)\}$ is also a basis for $\mathbb{R}^3$.

Example 20.13. Let $\mathbb{Q}(\sqrt{2}) = \{a + b\sqrt{2} : a, b \in \mathbb{Q}\}$. The sets $\{1, \sqrt{2}\}$ and $\{1 + \sqrt{2}, 1 - \sqrt{2}\}$ are both bases of $\mathbb{Q}(\sqrt{2})$.

From the last two examples it should be clear that a given vector space has several bases. In fact, there are an infinite number of bases for both of these examples. *In general, there is no unique basis for a vector space.* However, every basis of $\mathbb{R}^3$ consists of exactly three vectors, and every basis of $\mathbb{Q}(\sqrt{2})$ consists of exactly two vectors. This is a consequence of the next proposition.

Proposition 20.14. *Let $\{e_1, e_2, \ldots, e_m\}$ and $\{f_1, f_2, \ldots, f_n\}$ be two bases for a vector space V. Then $m = n$.*

PROOF.   Since $\{e_1, e_2, \ldots, e_m\}$ is a basis, it is a linearly independent set. By Proposition 20.11, $n \leq m$. Similarly, $\{f_1, f_2, \ldots, f_n\}$ is a linearly independent set, and the last proposition implies that $m \leq n$. Consequently, $m = n$.   □

If $\{e_1, e_2, \ldots, e_n\}$ is a basis for a vector space V, then we say that the ***dimension*** of V is n and we write $\dim V = n$. We will leave the proof of the following theorem as an exercise.

Theorem 20.15. *Let V be a vector space of dimension n.*

1. *If $S = \{v_1, \ldots, v_n\}$ is a set of linearly independent vectors for V, then S is a basis for V.*

2. *If $S = \{v_1, \ldots, v_n\}$ spans V, then S is a basis for V.*

3. *If $S = \{v_1, \ldots, v_k\}$ is a set of linearly independent vectors for V with $k < n$, then there exist vectors $v_{k+1}, \ldots, v_n$ such that*

$$\{v_1, \ldots, v_k, v_{k+1}, \ldots, v_n\}$$

is a basis for V.

20.4   Exercises

1.   If F is a field, show that $F[x]$ is a vector space over F, where the vectors in $F[x]$ are polynomials. Vector addition is polynomial addition, and scalar multiplication is defined by $\alpha p(x)$ for $\alpha \in F$.

2.   Prove that $\mathbb{Q}(\sqrt{2})$ is a vector space.

3.   Let $\mathbb{Q}(\sqrt{2}, \sqrt{3})$ be the field generated by elements of the form $a + b\sqrt{2} + c\sqrt{3} + d\sqrt{6}$, where a, b, c, d are in $\mathbb{Q}$. Prove that $\mathbb{Q}(\sqrt{2}, \sqrt{3})$ is a vector space of dimension 4 over $\mathbb{Q}$. Find a basis for $\mathbb{Q}(\sqrt{2}, \sqrt{3})$.

4.   Prove that the complex numbers are a vector space of dimension 2 over $\mathbb{R}$.

5.   Prove that the set P_n of all polynomials of degree less than n form a subspace of the vector space $F[x]$. Find a basis for P_n and compute the dimension of P_n.

6.   Let F be a field and denote the set of n-tuples of F by F^n. Given vectors $u = (u_1, \ldots, u_n)$ and $v = (v_1, \ldots, v_n)$ in F^n and α in F, define vector addition by

$$u + v = (u_1, \ldots, u_n) + (v_1, \ldots, v_n) = (u_1 + v_1, \ldots, u_n + v_n)$$

and scalar multiplication by

$$\alpha u = \alpha(u_1, \ldots, u_n) = (\alpha u_1, \ldots, \alpha u_n).$$

Prove that F^n is a vector space of dimension n under these operations.

7.   Which of the following sets are subspaces of $\mathbb{R}^3$? If the set is indeed a subspace, find a basis for the subspace and compute its dimension.

(a) $\{(x_1, x_2, x_3) : 3x_1 - 2x_2 + x_3 = 0\}$

(b) $\{(x_1, x_2, x_3) : 3x_1 + 4x_3 = 0, 2x_1 - x_2 + x_3 = 0\}$

(c) $\{(x_1, x_2, x_3) : x_1 - 2x_2 + 2x_3 = 2\}$

(d) $\{(x_1, x_2, x_3) : 3x_1 - 2x_2^2 = 0\}$

8. Show that the set of all possible solutions $(x, y, z) \in \mathbb{R}^3$ of the equations

$$Ax + By + Cz = 0$$
$$Dx + Ey + Cz = 0$$

form a subspace of $\mathbb{R}^3$.

9. Let W be the subset of continuous functions on $[0, 1]$ such that $f(0) = 0$. Prove that W is a subspace of $C[0, 1]$.

10. Let V be a vector space over F. Prove that $-(\alpha v) = (-\alpha)v = \alpha(-v)$ for all $\alpha \in F$ and all $v \in V$.

11. Let V be a vector space of dimension n. Prove each of the following statements.

(a) If $S = \{v_1, \ldots, v_n\}$ is a set of linearly independent vectors for V, then S is a basis for V.

(b) If $S = \{v_1, \ldots, v_n\}$ spans V, then S is a basis for V.

(c) If $S = \{v_1, \ldots, v_k\}$ is a set of linearly independent vectors for V with $k < n$, then there exist vectors $v_{k+1}, \ldots, v_n$ such that

$$\{v_1, \ldots, v_k, v_{k+1}, \ldots, v_n\}$$

is a basis for V.

12. Prove that any set of vectors containing $\mathbf{0}$ is linearly dependent.

13. Let V be a vector space. Show that $\{\mathbf{0}\}$ is a subspace of V of dimension zero.

14. If a vector space V is spanned by n vectors, show that any set of m vectors in V must be linearly dependent for $m > n$.

15. Linear Transformations. Let V and W be vector spaces over a field F, of dimensions m and n, respectively. If $T : V \to W$ is a map satisfying

$$T(u + v) = T(u) + T(v)$$
$$T(\alpha v) = \alpha T(v)$$

for all $\alpha \in F$ and all $u, v \in V$, then T is called a ***linear transformation*** from V into W.

(a) Prove that the ***kernel*** of T, $\ker(T) = \{v \in V : T(v) = \mathbf{0}\}$, is a subspace of V. The kernel of T is sometimes called the ***null space*** of T.

(b) Prove that the ***range*** or ***range space*** of T, $R(V) = \{w \in W : T(v) = w$ for some $v \in V\}$, is a subspace of W.

(c) Show that $T : V \to W$ is injective if and only if $\ker(T) = \{\mathbf{0}\}$.

(d) Let $\{v_1, \ldots, v_k\}$ be a basis for the null space of T. We can extend this basis to be a basis $\{v_1, \ldots, v_k, v_{k+1}, \ldots, v_m\}$ of V. Why? Prove that $\{T(v_{k+1}), \ldots, T(v_m)\}$ is a basis for the range of T. Conclude that the range of T has dimension $m - k$.

(e) Let $\dim V = \dim W$. Show that a linear transformation $T : V \to W$ is injective if and only if it is surjective.

16.   Let V and W be finite dimensional vector spaces of dimension n over a field F. Suppose that $T : V \to W$ is a vector space isomorphism. If $\{v_1, \ldots, v_n\}$ is a basis of V, show that $\{T(v_1), \ldots, T(v_n)\}$ is a basis of W. Conclude that any vector space over a field F of dimension n is isomorphic to F^n.

17. **Direct Sums.**   Let U and V be subspaces of a vector space W. The sum of U and V, denoted $U + V$, is defined to be the set of all vectors of the form $u + v$, where $u \in U$ and $v \in V$.

(a) Prove that $U + V$ and $U \cap V$ are subspaces of W.

(b) If $U + V = W$ and $U \cap V = \mathbf{0}$, then W is said to be the **direct sum.** In this case, we write $W = U \oplus V$. Show that every element $w \in W$ can be written uniquely as $w = u + v$, where $u \in U$ and $v \in V$.

(c) Let U be a subspace of dimension k of a vector space W of dimension n. Prove that there exists a subspace V of dimension $n - k$ such that $W = U \oplus V$. Is the subspace V unique?

(d) If U and V are arbitrary subspaces of a vector space W, show that

$$\dim(U + V) = \dim U + \dim V - \dim(U \cap V).$$

18. **Dual Spaces.**   Let V and W be finite dimensional vector spaces over a field F.

(a) Show that the set of all linear transformations from V into W, denoted $\operatorname{Hom}(V, W)$, is a vector space over F, where we define vector addition as follows:

$$(S + T)(v) = S(v) + T(v)$$
$$(\alpha S)(v) = \alpha S(v),$$

where $S, T \in \operatorname{Hom}(V, W)$, $\alpha \in F$, and $v \in V$.

(b) Let V be an F-vector space. Define the **dual space** of V to be $V^* = \operatorname{Hom}(V, F)$. Elements in the dual space of V are called **linear functionals.** Let $v_1, \ldots, v_n$ be an ordered basis for V. If $v = \alpha_1 v_1 + \cdots + \alpha_n v_n$ is any vector in V, define a linear functional $\phi_i : V \to F$ by $\phi_i(v) = \alpha_i$. Show that the ϕ_i's form a basis for V^*. This basis is called the **dual basis** of $v_1, \ldots, v_n$ (or simply the dual basis if the context makes the meaning clear).

(c) Consider the basis $\{(3, 1), (2, -2)\}$ for $\mathbb{R}^2$. What is the dual basis for $(\mathbb{R}^2)^*$?

(d) Let V be a vector space of dimension n over a field F and let V^{**} be the dual space of V^*. Show that each element $v \in V$ gives rise to an element λ_v in V^{**} and that the map $v \mapsto \lambda_v$ is an isomorphism of V with V^{**}.

References and Suggested Readings

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[5]   Johnson, L. W., Riess, R. D., and Arnold, J. T. *Introduction to Linear Algebra.* 6th ed. Pearson, Upper Saddle River, NJ, 2011.

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20.5   Sage

Many computations, in seemingly very different areas of mathematics, can be translated into questions about linear combinations, or other areas of linear algebra. So Sage has extensive and thorough support for topics such as vector spaces.

Vector Spaces

The simplest way to create a vector space is to begin with a field and use an exponent to indicate the number of entries in the vectors of the space.

```
V = QQ^4; V
```

Vector space of dimension 4 over Rational Field

```
F.<a> = FiniteField(3^4)
W = F^5; W
```

Vector space of dimension 5 over Finite Field **in** a of size 3^4

Elements can be built with the vector constructor.

```
v = vector(QQ, [1, 1/2, 1/3, 1/4]); v
```

(1, 1/2, 1/3, 1/4)

```
v in V
```

True

```
w = vector(F, [1, a^2, a^4, a^6, a^8]); w
```

(1, a^2, a^3 + 1, a^3 + a^2 + a + 1, a^2 + a + 2)

```
w in W
```

True

Notice that vectors are printed with parentheses, which helps distinguish them from lists (though they alos look like tuples). Vectors print horizontally, but in Sage there is no such thing as a "row vector" or a "column vector," though once matrices get involved we need to address this distinction. Finally, notice how the elements of the finite field have been converted to an alternate representation.

Once we have vector spaces full of vectors, we can perform computations with them. Ultimately, all the action in a vector space comes back to vector addition and scalar multiplication, which together create linear combinations.

```
u = vector(QQ, [ 1, 2,  3, 4,   5,   6])
v = vector(QQ, [-1, 2, -4, 8, -16, 32])
3*u - 2*v
```

```
(5, 2, 17, -4, 47, -46)
```

```
w = vector(F, [1, a^2, a^4, a^6,  a^8])
x = vector(F, [1,   a, 2*a,   a,    1])
y = vector(F, [1, a^3, a^6, a^9, a^12])
a^25*w + a^43*x + a^66*y
```

```
(a^3 + a^2 + a + 2, a^2 + 2*a, 2*a^3 + a^2 + 2, 2*a^3 + a^2 + a,
 a^3 + 2*a^2 + a + 2)
```

Subspaces

Sage can create subspaces in a variety of ways, such as in the creation of row or column spaces of matrices. However, the most direct way is to begin with a set of vectors to use as a spanning set.

```
u = vector(QQ, [1, -1, 3])
v = vector(QQ, [2, 1, -1])
w = vector(QQ, [3, 0, 2])
S = (QQ^3).subspace([u, v, w]); S
```

```
Vector space of degree 3 and dimension 2 over Rational Field
Basis matrix:
[   1    0  2/3]
[   0    1 -7/3]
```

```
3*u - 6*v + (1/2)*w in S
```

```
True
```

```
vector(QQ, [4, -1, -2]) in S
```

```
False
```

Notice that the information printed about S includes a "basis matrix." The rows of this matrix are a basis for the vector space. We can get the basis, as a list of vectors (not rows of a matrix), with the `.basis()` method.

```
S.basis()
```

```
[
(1, 0, 2/3),
(0, 1, -7/3)
]
```

Notice that Sage has converted the spanning set of three vectors into a basis with two vectors. This is partially due to the fact that the original set of three vectors is linearly dependent, but a more substantial change has occurred.

This is a good place to discuss some of the mathematics behind what makes Sage work. A vector space over an infinite field, like the rationals or the reals, is an infinite set. No matter how expansive computer memory may seem, it is still finite. How does Sage fit an infinite set into our finite machines? The main idea is that a finite-dimensional vector

space has a finite set of generators, which we know as a basis. So Sage really only needs the elements of a basis (two vectors in the previous example) to be able to work with the infinitely many possibilities for elements of the subspace.

Furthermore, for every basis associated with a vector space, Sage performs linear combinations to convert the given basis into another "standard" basis. This new basis has the property that as the rows of a matrix, the matrix is in reduced row-echelon form. You can see this in the basis matrix above. The reduced row-echelon form of a matrix is unique, so this standard basis allows Sage to recognize when two vector spaces are equal. Here is an example.

```
u = vector(QQ, [1, -1,  3])
v = vector(QQ, [2,  1, -1])
w = vector(QQ, [3,  0,  2])
u + v == w
```

True

```
S1 = (QQ^3).subspace([u, v, w])
S2 = (QQ^3).subspace([u-v, v-w, w-u])
S1 == S2
```

True

As you might expect, it is easy to determine the dimension of a vector space.

```
u = vector(QQ, [1, -1,  3,  4])
v = vector(QQ, [2,  1, -1, -2])
S = (QQ^4).subspace([u, v, 2*u + 3*v, -u + 2*v])
S.dimension()
```

2

Linear Independence

There are a variety of ways in Sage to determine if a set of vectors is linearly independent or not, and to find relations of linear dependence if they exist. The technique we will show here is a simple test to see if a set of vectors is linearly independent or not. Simply use the vectors as a spanning set for a subspace, and check the dimension of the subspace. The dimension equals the number of vectors in the spanning set if and only if the spanning set is linearly independent.

```
F.<a> = FiniteField(3^4)
u = vector(F, [a^i for i in range(0,  7, 1)])
v = vector(F, [a^i for i in range(0, 14, 2)])
w = vector(F, [a^i for i in range(0, 21, 3)])
S = (F^7).subspace([u, v, w])
S.dimension()
```

3

```
S = (F^7).subspace([u, v, a^3*u + a^11*v])
S.dimension()
```

2

So the first set of vectors, [u, v, w], is linearly independent, while the second set, [u, v, a^3*u + a^11*v], is not.

Abstract Vector Spaces

Sage does not implement many abstract vector spaces directly, such as P_n, the vector space of polynomials of degree n or less. This is due in part to the fact that a finite-dimensional vector space over a field F is isomorphic to the vector space F^n. So Sage captures all the functionality of finite-dimensional vector spaces, and it is left to the user to perform the conversions according to the isomorphism (which is often trivial with the choice of an obvious basis).

However, there are instances where rings behave naturally as vector spaces and we can exploit this extra structure. We will see much more of this in the chapters on fields and Galois theory. As an example, finite fields have a single generator, and the first few powers of the generator form a basis. Consider creating a vector space from the elements of a finite field of order $7^6 = 117\,649$. As elements of a field we know they can be added, so we will *define* this to be the addition in our vector space. For any element of the integers mod 7, we can multiply an element of the field by the integer, so we *define* this to be our scalar multiplication. Later, we will be certain that these two definitions lead to a vector space, but take that for granted now. So here are some operations in our new vector space.

```
F.<a> = FiniteField(7^6)
u = 2*a^5 + 6*a^4 + 2*a^3 + 3*a^2 + 2*a + 3
v = 4*a^5 + 4*a^4 + 4*a^3 + 6*a^2 + 5*a + 6
u + v
```

```
6*a^5 + 3*a^4 + 6*a^3 + 2*a^2 + 2
```

```
4*u
```

```
a^5 + 3*a^4 + a^3 + 5*a^2 + a + 5
```

```
2*u + 5*v
```

```
3*a^5 + 4*a^4 + 3*a^3 + a^2 + a + 1
```

You might recognize that this looks very familiar to how we add polynomials, and multiply polynomials by scalars. You would be correct. However, notice that in this vector space construction, we are totally ignoring the possibility of multiplying two field elements together. As a vector space with scalars from $\mathbb{Z}_7$, a basis is the first six powers of the generator, $\{1, a, a^2, a^3, a^4, a^5\}$. (Notice how counting from zero is natural here.) You may have noticed how Sage consistently rewrites elements of fields as linear combinations — now you have a good explanation.

Here is what Sage knows about a finite field as a vector space. First, it knows that the finite field *is* a vector space, and what the field of scalars is.

```
V = F.vector_space(); V
```

```
Vector space of dimension 6 over Finite Field of size 7
```

```
R = V.base_ring(); R
```

```
Finite Field of size 7
```

```
R == FiniteField(7)
```

```
True
```

```
V.dimension()
```

6

So the finite field (as a vector space) is isomorphic to the vector space $(\mathbb{Z}_7)^6$. Notice this is not a ring or field isomorphism, as it does not fully address multiplication of elements, even though that is possible in the field.

Second, elements of the field can be converted to elements of the vector space easily.

```
x = V(u); x
```

(3, 2, 3, 2, 6, 2)

```
y = V(v); y
```

(6, 5, 6, 4, 4, 4)

Notice that Sage writes field elements with high powers of the generator first, while the basis in use is ordered with low powers first. The computations below illustrate the isomorphism preserving the structure between the finite field itself and its interpretation as a vector space, $(\mathbb{Z}_7)^6$.

```
V(u + v) == V(u) + V(v)
```

True

```
two = R(2)
V(two*u) == two*V(u)
```

True

Linear Algebra

Sage has extensive support for linear algebra, well beyond what we have described here, or what we will need for the remaining chapters. Create vector spaces and vectors (with different fields of scalars), and then use tab-completion on these objects to explore the large sets of available commands.

20.6 Sage Exercises

1. Given two subspaces U and W of a vector space V, their sum $U + W$ can be defined as the *set* $U + W = \{u + w \mid u \in U, \ w \in W\}$, in other words, the set of all possible sums of an element from U and an element from W.

Notice this is not the direct sum of your text, nor the `direct_sum()` method in Sage. However, you can build this subspace in Sage as follows. Grab the bases of U and W individually, as lists of vectors. Join the two lists together by just using a plus sign between them. Now build the sum subspace by creating a subspace of V spanned by this set, by using the `.subspace()` method.

In the vector space (`QQ^10`) construct two subspaces that you expect to (a) have dimension 5 or 6 or so, and (b) have an intersection that is a vector space of dimension 2 or so. Compare their individual dimensions with the dimensions of the intersection of U and W ($U \cap W$, `.intersection()` in Sage) and the sum $U + W$.

Repeat the experiment with the two original vector spaces having dimension 8 or so, and with the intersection as small as possible. Form a general conjecture relating these four dimensions based on the results of your two (or more)experiments.

2.   We can construct a field in Sage that extends the rationals by adding in a fourth root of two, $\mathbb{Q}[\sqrt[4]{2}]$, with the command `F.<c> = QQ[2^(1/4)]`. This is a vector space of dimension 4 over the rationals, with a basis that is the first four powers of $c = \sqrt[4]{2}$ (starting with the zero power).

The command `F.vector_space()` will return three items in a triple (so be careful how you handle this output to extract what you need). The first part of the output is a vector space over the rationals that is isomorphic to `F`. The next is a vector space isomorphism (invertible linear transformation) from the provided vector space to the field, while the third is an isomorphism in the opposite direction. These two isomorphisms can then be used like functions. Notice that this is different behavior than for `.vector_space()` applied to finite fields. Create non-trivial examples that show that these vector space isomorphisms behave as an isomorphism should. (You will have at least four such examples in a complete solution.)

3.   Build a finite field F of order p^n in the usual way. Then construct the (multiplicative) group of all invertible (nonsingular) $m \times m$ matrices over this field with the command `G = GL(m, F)` ("the general linear group"). What is the order of this group? In other words, find a general expression for the order of this group.

Your answer should be a function of m, p and n. Provide a complete explanation of the logic behind your solution (i.e. something resembling a proof). Also provide tests in Sage that your answer is correct.

Hints: `G.order()` will help you test and verify your hypotheses. Small examples in Sage (listing all the elements of the group) might aid your intuition—which is why this is a Sage exercise. Small means 2×2 and 3×3 matrices and finite fields with $2, 3, 4, 5$ elements, at most. Results do not really depend on each of p and n, but rather just on p^n.

Realize this group is interesting because it contains representations of all the invertible (i.e. 1-1 and onto) linear transformations from the (finite) vector space F^m to itself.

4.   What happens if we try to do linear algebra over a *ring* that is not also a *field*? The object that resembles a vector space, but with this one distinction, is known as a **module**. You can build one easily with a construction like `ZZ^3`. Evaluate the following to create a module and a submodule.

```
M = ZZ^3
u = M([1, 0, 0])
v = M([2, 2, 0])
w = M([0, 0, 4])
N = M.submodule([u, v, w])
```

Examine the bases and dimensions (aka "rank") of the module and submodule, and check the equality of the module and submodule. How is this different than the situation for vector spaces? Can you create a third module, `P`, that is a proper subset of `M` and properly contains `N`?

5.   A finite field, F, of order 5^3 is a vector space of dimension 3 over $\mathbb{Z}_5$. Suppose a is a generator of F. Let M be any 3×3 matrix with entries from $\mathbb{Z}_5$ (carefule here, the elements are from th field of scalars, not from the vector space). If we convert an element $x \in F$ to a vector (relative to the basis $\{1, a, a^2\}$), then we can multiply it by M (with M on the left) to create another vector, which we can translate to a linear combination of the basis elements, and hence another element of F. This function is a vector space homomorphism, better known as a linear transformation (implemented with a matrix representation relative to the basis $\{1, a, a^2\}$. Notice that each part below becomes less general and more specific.

(a) Create a non-invertible matrix R and give examples to show that the mapping described by R is a vector space homomorphism of F into F.

(b) Create an invertible matrix M. The mapping will now be an invertible homomorphism. Determine the inverse function and give examples to verify its properties.

(c) Since a is a generator of the field, the mapping $a \mapsto a^5$ can be extended to a vector space homomorphism (i.e. a linear transformation). Find a matrix M which effects this linear transformation, and from this, determine that the homomorphism is invertible.

(d) None of the previous three parts applies to properties of multiplication in the field. However, the mapping from the third part also preserves multiplication in the field, though a proof of this may not be obvious right now. So we are saying this mapping is a field automorphism, preserving both addition and multiplication. Give a nontrivial example of the multiplication-preserving properties of this mapping. (This is the **Frobenius map** which will be discussed further in Chapter 21.)

Fields

It is natural to ask whether or not some field F is contained in a larger field. We think of the rational numbers, which reside inside the real numbers, while in turn, the real numbers live inside the complex numbers. We can also study the fields between $\mathbb{Q}$ and $\mathbb{R}$ and inquire as to the nature of these fields.

More specifically if we are given a field F and a polynomial $p(x) \in F[x]$, we can ask whether or not we can find a field E containing F such that $p(x)$ factors into linear factors over $E[x]$. For example, if we consider the polynomial

$$p(x) = x^4 - 5x^2 + 6$$

in $\mathbb{Q}[x]$, then $p(x)$ factors as $(x^2 - 2)(x^2 - 3)$. However, both of these factors are irreducible in $\mathbb{Q}[x]$. If we wish to find a zero of $p(x)$, we must go to a larger field. Certainly the field of real numbers will work, since

$$p(x) = (x - \sqrt{2})(x + \sqrt{2})(x - \sqrt{3})(x + \sqrt{3}).$$

It is possible to find a smaller field in which $p(x)$ has a zero, namely

$$\mathbb{Q}(\sqrt{2}) = \{a + b\sqrt{2} : a, b \in \mathbb{Q}\}.$$

We wish to be able to compute and study such fields for arbitrary polynomials over a field F.

21.1 Extension Fields

A field E is an **extension field** of a field F if F is a subfield of E. The field F is called the **base field**. We write $F \subset E$.

Example 21.1. For example, let

$$F = \mathbb{Q}(\sqrt{2}) = \{a + b\sqrt{2} : a, b \in \mathbb{Q}\}$$

and let $E = \mathbb{Q}(\sqrt{2} + \sqrt{3})$ be the smallest field containing both $\mathbb{Q}$ and $\sqrt{2} + \sqrt{3}$. Both E and F are extension fields of the rational numbers. We claim that E is an extension field of F. To see this, we need only show that $\sqrt{2}$ is in E. Since $\sqrt{2} + \sqrt{3}$ is in E, $1/(\sqrt{2} + \sqrt{3}) = \sqrt{3} - \sqrt{2}$ must also be in E. Taking linear combinations of $\sqrt{2} + \sqrt{3}$ and $\sqrt{3} - \sqrt{2}$, we find that $\sqrt{2}$ and $\sqrt{3}$ must both be in E.

Example 21.2. Let $p(x) = x^2 + x + 1 \in \mathbb{Z}_2[x]$. Since neither 0 nor 1 is a root of this polynomial, we know that $p(x)$ is irreducible over $\mathbb{Z}_2$. We will construct a field extension of $\mathbb{Z}_2$ containing an element α such that $p(\alpha) = 0$. By Theorem 17.22, the ideal $\langle p(x) \rangle$

generated by $p(x)$ is maximal; hence, $\mathbb{Z}_2[x]/\langle p(x)\rangle$ is a field. Let $f(x) + \langle p(x)\rangle$ be an arbitrary element of $\mathbb{Z}_2[x]/\langle p(x)\rangle$. By the division algorithm,

$$f(x) = (x^2 + x + 1)q(x) + r(x),$$

where the degree of $r(x)$ is less than the degree of $x^2 + x + 1$. Therefore,

$$f(x) + \langle x^2 + x + 1\rangle = r(x) + \langle x^2 + x + 1\rangle.$$

The only possibilities for $r(x)$ are then 0, 1, x, and $1+x$. Consequently, $E = \mathbb{Z}_2[x]/\langle x^2+x+1\rangle$ is a field with four elements and must be a field extension of $\mathbb{Z}_2$, containing a zero α of $p(x)$. The field $\mathbb{Z}_2(\alpha)$ consists of elements

$$0 + 0\alpha = 0$$
$$1 + 0\alpha = 1$$
$$0 + 1\alpha = \alpha$$
$$1 + 1\alpha = 1 + \alpha.$$

Notice that $\alpha^2 + \alpha + 1 = 0$; hence, if we compute $(1 + \alpha)^2$,

$$(1 + \alpha)(1 + \alpha) = 1 + \alpha + \alpha + (\alpha)^2 = \alpha.$$

Other calculations are accomplished in a similar manner. We summarize these computations in the following tables, which tell us how to add and multiply elements in E.

$+$	0	1	α	$1+\alpha$
0	0	1	α	$1+\alpha$
1	1	0	$1+\alpha$	α
α	α	$1+\alpha$	0	1
$1+\alpha$	$1+\alpha$	α	1	0

Table 21.3: Addition Table for $\mathbb{Z}_2(\alpha)$

$\cdot$	0	1	α	$1+\alpha$
0	0	0	0	0
1	0	1	α	$1+\alpha$
α	0	α	$1+\alpha$	1
$1+\alpha$	0	$1+\alpha$	1	α

Table 21.4: Multiplication Table for $\mathbb{Z}_2(\alpha)$

The following theorem, due to Kronecker, is so important and so basic to our understanding of fields that it is often known as the Fundamental Theorem of Field Theory.

Theorem 21.5. *Let F be a field and let $p(x)$ be a nonconstant polynomial in $F[x]$. Then there exists an extension field E of F and an element $\alpha \in E$ such that $p(\alpha) = 0$.*

PROOF.   To prove this theorem, we will employ the method that we used to construct Example 21.2. Clearly, we can assume that $p(x)$ is an irreducible polynomial. We wish to find an extension field E of F containing an element α such that $p(\alpha) = 0$. The ideal $\langle p(x) \rangle$ generated by $p(x)$ is a maximal ideal in $F[x]$ by Theorem 17.22; hence, $F[x]/\langle p(x) \rangle$ is a field. We claim that $E = F[x]/\langle p(x) \rangle$ is the desired field.

We first show that E is a field extension of F. We can define a homomorphism of commutative rings by the map $\psi : F \to F[x]/\langle p(x) \rangle$, where $\psi(a) = a + \langle p(x) \rangle$ for $a \in F$. It is easy to check that ψ is indeed a ring homomorphism. Observe that

$$\psi(a) + \psi(b) = (a + \langle p(x) \rangle) + (b + \langle p(x) \rangle) = (a + b) + \langle p(x) \rangle = \psi(a + b)$$

and

$$\psi(a)\psi(b) = (a + \langle p(x) \rangle)(b + \langle p(x) \rangle) = ab + \langle p(x) \rangle = \psi(ab).$$

To prove that ψ is one-to-one, assume that

$$a + \langle p(x) \rangle = \psi(a) = \psi(b) = b + \langle p(x) \rangle.$$

Then $a - b$ is a multiple of $p(x)$, since it lives in the ideal $\langle p(x) \rangle$. Since $p(x)$ is a nonconstant polynomial, the only possibility is that $a - b = 0$. Consequently, $a = b$ and ψ is injective. Since ψ is one-to-one, we can identify F with the subfield $\{a + \langle p(x) \rangle : a \in F\}$ of E and view E as an extension field of F.

It remains for us to prove that $p(x)$ has a zero $\alpha \in E$. Set $\alpha = x + \langle p(x) \rangle$. Then α is in E. If $p(x) = a_0 + a_1 x + \cdots + a_n x^n$, then

$$\begin{aligned} p(\alpha) &= a_0 + a_1(x + \langle p(x) \rangle) + \cdots + a_n(x + \langle p(x) \rangle)^n \\ &= a_0 + (a_1 x + \langle p(x) \rangle) + \cdots + (a_n x^n + \langle p(x) \rangle) \\ &= a_0 + a_1 x + \cdots + a_n x^n + \langle p(x) \rangle \\ &= 0 + \langle p(x) \rangle. \end{aligned}$$

Therefore, we have found an element $\alpha \in E = F[x]/\langle p(x) \rangle$ such that α is a zero of $p(x)$.   □

Example 21.6. Let $p(x) = x^5 + x^4 + 1 \in \mathbb{Z}_2[x]$. Then $p(x)$ has irreducible factors $x^2 + x + 1$ and $x^3 + x + 1$. For a field extension E of $\mathbb{Z}_2$ such that $p(x)$ has a root in E, we can let E be either $\mathbb{Z}_2[x]/\langle x^2 + x + 1 \rangle$ or $\mathbb{Z}_2[x]/\langle x^3 + x + 1 \rangle$. We will leave it as an exercise to show that $\mathbb{Z}_2[x]/\langle x^3 + x + 1 \rangle$ is a field with $2^3 = 8$ elements.

Algebraic Elements

An element α in an extension field E over F is **algebraic** over F if $f(\alpha) = 0$ for some nonzero polynomial $f(x) \in F[x]$. An element in E that is not algebraic over F is **transcendental** over F. An extension field E of a field F is an **algebraic extension** of F if every element in E is algebraic over F. If E is a field extension of F and $\alpha_1, \ldots, \alpha_n$ are contained in E, we denote the smallest field containing F and $\alpha_1, \ldots, \alpha_n$ by $F(\alpha_1, \ldots, \alpha_n)$. If $E = F(\alpha)$ for some $\alpha \in E$, then E is a **simple extension** of F.

Example 21.7. Both $\sqrt{2}$ and i are algebraic over $\mathbb{Q}$ since they are zeros of the polynomials $x^2 - 2$ and $x^2 + 1$, respectively. Clearly π and e are algebraic over the real numbers; however, it is a nontrivial fact that they are transcendental over $\mathbb{Q}$. Numbers in $\mathbb{R}$ that are algebraic over $\mathbb{Q}$ are in fact quite rare. Almost all real numbers are transcendental over $\mathbb{Q}$.[1] (In many cases we do not know whether or not a particular number is transcendental; for example, it is still not known whether $\pi + e$ is transcendental or algebraic.)

---

[1] The probability that a real number chosen at random from the interval $[0, 1]$ will be transcendental over the rational numbers is one.

A complex number that is algebraic over $\mathbb{Q}$ is an **algebraic number**. A **transcendental number** is an element of $\mathbb{C}$ that is transcendental over $\mathbb{Q}$.

Example 21.8. We will show that $\sqrt{2 + \sqrt{3}}$ is algebraic over $\mathbb{Q}$. If $\alpha = \sqrt{2 + \sqrt{3}}$, then $\alpha^2 = 2 + \sqrt{3}$. Hence, $\alpha^2 - 2 = \sqrt{3}$ and $(\alpha^2 - 2)^2 = 3$. Since $\alpha^4 - 4\alpha^2 + 1 = 0$, it must be true that α is a zero of the polynomial $x^4 - 4x^2 + 1 \in \mathbb{Q}[x]$.

It is very easy to give an example of an extension field E over a field F, where E contains an element transcendental over F. The following theorem characterizes transcendental extensions.

Theorem 21.9. *Let E be an extension field of F and $\alpha \in E$. Then α is transcendental over F if and only if $F(\alpha)$ is isomorphic to $F(x)$, the field of fractions of $F[x]$.*

PROOF. Let $\phi_\alpha : F[x] \to E$ be the evaluation homomorphism for α. Then α is transcendental over F if and only if $\phi_\alpha(p(x)) = p(\alpha) \neq 0$ for all nonconstant polynomials $p(x) \in F[x]$. This is true if and only if $\ker \phi_\alpha = \{0\}$; that is, it is true exactly when ϕ_α is one-to-one. Hence, E must contain a copy of $F[x]$. The smallest field containing $F[x]$ is the field of fractions $F(x)$. By Theorem 18.4, E must contain a copy of this field. $\square$

We have a more interesting situation in the case of algebraic extensions.

Theorem 21.10. *Let E be an extension field of a field F and $\alpha \in E$ with α algebraic over F. Then there is a unique irreducible monic polynomial $p(x) \in F[x]$ of smallest degree such that $p(\alpha) = 0$. If $f(x)$ is another polynomial in $F[x]$ such that $f(\alpha) = 0$, then $p(x)$ divides $f(x)$.*

PROOF. Let $\phi_\alpha : F[x] \to E$ be the evaluation homomorphism. The kernel of ϕ_α is a principal ideal generated by some $p(x) \in F[x]$ with $\deg p(x) \geq 1$. We know that such a polynomial exists, since $F[x]$ is a principal ideal domain and α is algebraic. The ideal $\langle p(x) \rangle$ consists exactly of those elements of $F[x]$ having α as a zero. If $f(\alpha) = 0$ and $f(x)$ is not the zero polynomial, then $f(x) \in \langle p(x) \rangle$ and $p(x)$ divides $f(x)$. So $p(x)$ is a polynomial of minimal degree having α as a zero. Any other polynomial of the same degree having α as a zero must have the form $\beta p(x)$ for some $\beta \in F$.

Suppose now that $p(x) = r(x)s(x)$ is a factorization of $p(x)$ into polynomials of lower degree. Since $p(\alpha) = 0$, $r(\alpha)s(\alpha) = 0$; consequently, either $r(\alpha) = 0$ or $s(\alpha) = 0$, which contradicts the fact that p is of minimal degree. Therefore, $p(x)$ must be irreducible. $\square$

Let E be an extension field of F and $\alpha \in E$ be algebraic over F. The unique monic polynomial $p(x)$ of the last theorem is called the **minimal polynomial** for α over F. The degree of $p(x)$ is the **degree of α over F**.

Example 21.11. Let $f(x) = x^2 - 2$ and $g(x) = x^4 - 4x^2 + 1$. These polynomials are the minimal polynomials of $\sqrt{2}$ and $\sqrt{2 + \sqrt{3}}$, respectively.

Proposition 21.12. *Let E be a field extension of F and $\alpha \in E$ be algebraic over F. Then $F(\alpha) \cong F[x]/\langle p(x) \rangle$, where $p(x)$ is the minimal polynomial of α over F.*

PROOF. Let $\phi_\alpha : F[x] \to E$ be the evaluation homomorphism. The kernel of this map is $\langle p(x) \rangle$, where $p(x)$ is the minimal polynomial of α. By the First Isomorphism Theorem for rings, the image of ϕ_α in E is isomorphic to $F(\alpha)$ since it contains both F and α. $\square$

Theorem 21.13. *Let $E = F(\alpha)$ be a simple extension of F, where $\alpha \in E$ is algebraic over F. Suppose that the degree of α over F is n. Then every element $\beta \in E$ can be expressed uniquely in the form*

$$\beta = b_0 + b_1\alpha + \cdots + b_{n-1}\alpha^{n-1}$$

for $b_i \in F$.

PROOF.   Since $\phi_\alpha(F[x]) \cong F(\alpha)$, every element in $E = F(\alpha)$ must be of the form $\phi_\alpha(f(x)) = f(\alpha)$, where $f(\alpha)$ is a polynomial in α with coefficients in F. Let

$$p(x) = x^n + a_{n-1}x^{n-1} + \cdots + a_0$$

be the minimal polynomial of α. Then $p(\alpha) = 0$; hence,

$$\alpha^n = -a_{n-1}\alpha^{n-1} - \cdots - a_0.$$

Similarly,

$$
\begin{aligned}
\alpha^{n+1} &= \alpha\alpha^n \\
&= -a_{n-1}\alpha^n - a_{n-2}\alpha^{n-1} - \cdots - a_0\alpha \\
&= -a_{n-1}(-a_{n-1}\alpha^{n-1} - \cdots - a_0) - a_{n-2}\alpha^{n-1} - \cdots - a_0\alpha.
\end{aligned}
$$

Continuing in this manner, we can express every monomial α^m, $m \geq n$, as a linear combination of powers of α that are less than n. Hence, any $\beta \in F(\alpha)$ can be written as

$$\beta = b_0 + b_1\alpha + \cdots + b_{n-1}\alpha^{n-1}.$$

To show uniqueness, suppose that

$$\beta = b_0 + b_1\alpha + \cdots + b_{n-1}\alpha^{n-1} = c_0 + c_1\alpha + \cdots + c_{n-1}\alpha^{n-1}$$

for b_i and c_i in F. Then

$$g(x) = (b_0 - c_0) + (b_1 - c_1)x + \cdots + (b_{n-1} - c_{n-1})x^{n-1}$$

is in $F[x]$ and $g(\alpha) = 0$. Since the degree of $g(x)$ is less than the degree of $p(x)$, the irreducible polynomial of α, $g(x)$ must be the zero polynomial. Consequently,

$$b_0 - c_0 = b_1 - c_1 = \cdots = b_{n-1} - c_{n-1} = 0,$$

or $b_i = c_i$ for $i = 0, 1, \ldots, n-1$. Therefore, we have shown uniqueness. $\qquad\square$

Example 21.14. Since $x^2 + 1$ is irreducible over $\mathbb{R}$, $\langle x^2 + 1 \rangle$ is a maximal ideal in $\mathbb{R}[x]$. So $E = \mathbb{R}[x]/\langle x^2 + 1 \rangle$ is a field extension of $\mathbb{R}$ that contains a root of $x^2 + 1$. Let $\alpha = x + \langle x^2 + 1 \rangle$. We can identify E with the complex numbers. By Proposition 21.12, E is isomorphic to $\mathbb{R}(\alpha) = \{a + b\alpha : a, b \in \mathbb{R}\}$. We know that $\alpha^2 = -1$ in E, since

$$
\begin{aligned}
\alpha^2 + 1 &= (x + \langle x^2 + 1 \rangle)^2 + (1 + \langle x^2 + 1 \rangle) \\
&= (x^2 + 1) + \langle x^2 + 1 \rangle \\
&= 0.
\end{aligned}
$$

Hence, we have an isomorphism of $\mathbb{R}(\alpha)$ with $\mathbb{C}$ defined by the map that takes $a + b\alpha$ to $a + bi$.

Let E be a field extension of a field F. If we regard E as a vector space over F, then we can bring the machinery of linear algebra to bear on the problems that we will encounter in our study of fields. The elements in the field E are vectors; the elements in the field F are scalars. We can think of addition in E as adding vectors. When we multiply an element in E by an element of F, we are multiplying a vector by a scalar. This view of field extensions is especially fruitful if a field extension E of F is a finite dimensional vector space over F,

and Theorem 21.13 states that $E = F(\alpha)$ is finite dimensional vector space over F with basis $\{1, \alpha, \alpha^2, \ldots, \alpha^{n-1}\}$.

If an extension field E of a field F is a finite dimensional vector space over F of dimension n, then we say that E is a **_finite extension of degree_ n _over_ F**. We write

$$[E : F] = n.$$

to indicate the dimension of E over F.

Theorem 21.15. *Every finite extension field E of a field F is an algebraic extension.*

PROOF. Let $\alpha \in E$. Since $[E : F] = n$, the elements

$$1, \alpha, \ldots, \alpha^n$$

cannot be linearly independent. Hence, there exist $a_i \in F$, not all zero, such that

$$a_n \alpha^n + a_{n-1} \alpha^{n-1} + \cdots + a_1 \alpha + a_0 = 0.$$

Therefore,

$$p(x) = a_n x^n + \cdots + a_0 \in F[x]$$

is a nonzero polynomial with $p(\alpha) = 0$. $\square$

Remark 21.16. Theorem 21.15 says that every finite extension of a field F is an algebraic extension. The converse is false, however. We will leave it as an exercise to show that the set of all elements in $\mathbb{R}$ that are algebraic over $\mathbb{Q}$ forms an infinite field extension of $\mathbb{Q}$.

The next theorem is a counting theorem, similar to Lagrange's Theorem in group theory. Theorem 21.17 will prove to be an extremely useful tool in our investigation of finite field extensions.

Theorem 21.17. *If E is a finite extension of F and K is a finite extension of E, then K is a finite extension of F and*

$$[K : F] = [K : E][E : F].$$

PROOF. Let $\{\alpha_1, \ldots, \alpha_n\}$ be a basis for E as a vector space over F and $\{\beta_1, \ldots, \beta_m\}$ be a basis for K as a vector space over E. We claim that $\{\alpha_i \beta_j\}$ is a basis for K over F. We will first show that these vectors span K. Let $u \in K$. Then $u = \sum_{j=1}^{m} b_j \beta_j$ and $b_j = \sum_{i=1}^{n} a_{ij} \alpha_i$, where $b_j \in E$ and $a_{ij} \in F$. Then

$$u = \sum_{j=1}^{m} \left(\sum_{i=1}^{n} a_{ij} \alpha_i \right) \beta_j = \sum_{i,j} a_{ij} (\alpha_i \beta_j).$$

So the mn vectors $\alpha_i \beta_j$ must span K over F.

We must show that $\{\alpha_i \beta_j\}$ are linearly independent. Recall that a set of vectors $v_1, v_2, \ldots, v_n$ in a vector space V are linearly independent if

$$c_1 v_1 + c_2 v_2 + \cdots + c_n v_n = 0$$

implies that

$$c_1 = c_2 = \cdots = c_n = 0.$$

Let

$$u = \sum_{i,j} c_{ij} (\alpha_i \beta_j) = 0$$

for $c_{ij} \in F$. We need to prove that all of the c_{ij}'s are zero. We can rewrite u as

$$\sum_{j=1}^{m} \left(\sum_{i=1}^{n} c_{ij}\alpha_i \right) \beta_j = 0,$$

where $\sum_i c_{ij}\alpha_i \in E$. Since the β_j's are linearly independent over E, it must be the case that

$$\sum_{i=1}^{n} c_{ij}\alpha_i = 0$$

for all j. However, the α_j are also linearly independent over F. Therefore, $c_{ij} = 0$ for all i and j, which completes the proof.   □

The following corollary is easily proved using mathematical induction.

Corollary 21.18. *If F_i is a field for $i = 1, \ldots, k$ and F_{i+1} is a finite extension of F_i, then F_k is a finite extension of F_1 and*

$$[F_k : F_1] = [F_k : F_{k-1}] \cdots [F_2 : F_1].$$

Corollary 21.19. *Let E be an extension field of F. If $\alpha \in E$ is algebraic over F with minimal polynomial $p(x)$ and $\beta \in F(\alpha)$ with minimal polynomial $q(x)$, then $\deg q(x)$ divides $\deg p(x)$.*

PROOF.   We know that $\deg p(x) = [F(\alpha) : F]$ and $\deg q(x) = [F(\beta) : F]$. Since $F \subset F(\beta) \subset F(\alpha)$,

$$[F(\alpha) : F] = [F(\alpha) : F(\beta)][F(\beta) : F].   \qquad \square$$

Example 21.20. Let us determine an extension field of $\mathbb{Q}$ containing $\sqrt{3} + \sqrt{5}$. It is easy to determine that the minimal polynomial of $\sqrt{3} + \sqrt{5}$ is $x^4 - 16x^2 + 4$. It follows that

$$[\mathbb{Q}(\sqrt{3} + \sqrt{5}) : \mathbb{Q}] = 4.$$

We know that $\{1, \sqrt{3}\}$ is a basis for $\mathbb{Q}(\sqrt{3})$ over $\mathbb{Q}$. Hence, $\sqrt{3} + \sqrt{5}$ cannot be in $\mathbb{Q}(\sqrt{3})$. It follows that $\sqrt{5}$ cannot be in $\mathbb{Q}(\sqrt{3})$ either. Therefore, $\{1, \sqrt{5}\}$ is a basis for $\mathbb{Q}(\sqrt{3}, \sqrt{5}) = (\mathbb{Q}(\sqrt{3}))(\sqrt{5})$ over $\mathbb{Q}(\sqrt{3})$ and $\{1, \sqrt{3}, \sqrt{5}, \sqrt{3}\sqrt{5} = \sqrt{15}\}$ is a basis for $\mathbb{Q}(\sqrt{3}, \sqrt{5}) = \mathbb{Q}(\sqrt{3} + \sqrt{5})$ over $\mathbb{Q}$. This example shows that it is possible that some extension $F(\alpha_1, \ldots, \alpha_n)$ is actually a simple extension of F even though $n > 1$.

Example 21.21. Let us compute a basis for $\mathbb{Q}(\sqrt[3]{5}, \sqrt{5}\,i)$, where $\sqrt{5}$ is the positive square root of 5 and $\sqrt[3]{5}$ is the real cube root of 5. We know that $\sqrt{5}\,i \notin \mathbb{Q}(\sqrt[3]{5})$, so

$$[\mathbb{Q}(\sqrt[3]{5}, \sqrt{5}\,i) : \mathbb{Q}(\sqrt[3]{5})] = 2.$$

It is easy to determine that $\{1, \sqrt{5}i\}$ is a basis for $\mathbb{Q}(\sqrt[3]{5}, \sqrt{5}\,i)$ over $\mathbb{Q}(\sqrt[3]{5})$. We also know that $\{1, \sqrt[3]{5}, (\sqrt[3]{5})^2\}$ is a basis for $\mathbb{Q}(\sqrt[3]{5})$ over $\mathbb{Q}$. Hence, a basis for $\mathbb{Q}(\sqrt[3]{5}, \sqrt{5}\,i)$ over $\mathbb{Q}$ is

$$\{1, \sqrt{5}\,i, \sqrt[3]{5}, (\sqrt[3]{5})^2, (\sqrt[6]{5})^5 i, (\sqrt[6]{5})^7 i = 5\sqrt[6]{5}\,i \text{ or } \sqrt[6]{5}\,i\}.$$

Notice that $\sqrt[6]{5}\,i$ is a zero of $x^6 + 5$. We can show that this polynomial is irreducible over $\mathbb{Q}$ using Eisenstein's Criterion, where we let $p = 5$. Consequently,

$$\mathbb{Q} \subset \mathbb{Q}(\sqrt[6]{5}\,i) \subset \mathbb{Q}(\sqrt[3]{5}, \sqrt{5}\,i).$$

But it must be the case that $\mathbb{Q}(\sqrt[6]{5}\,i) = \mathbb{Q}(\sqrt[3]{5}, \sqrt{5}\,i)$, since the degree of both of these extensions is 6.

Theorem 21.22. *Let E be a field extension of F. Then the following statements are equivalent.*

1. *E is a finite extension of F.*

2. *There exists a finite number of algebraic elements $\alpha_1, \ldots, \alpha_n \in E$ such that $E = F(\alpha_1, \ldots, \alpha_n)$.*

3. *There exists a sequence of fields*

$$E = F(\alpha_1, \ldots, \alpha_n) \supset F(\alpha_1, \ldots, \alpha_{n-1}) \supset \cdots \supset F(\alpha_1) \supset F,$$

   where each field $F(\alpha_1, \ldots, \alpha_i)$ is algebraic over $F(\alpha_1, \ldots, \alpha_{i-1})$.

PROOF. $(1) \Rightarrow (2)$. Let E be a finite algebraic extension of F. Then E is a finite dimensional vector space over F and there exists a basis consisting of elements $\alpha_1, \ldots, \alpha_n$ in E such that $E = F(\alpha_1, \ldots, \alpha_n)$. Each α_i is algebraic over F by Theorem 21.15.

$(2) \Rightarrow (3)$. Suppose that $E = F(\alpha_1, \ldots, \alpha_n)$, where every α_i is algebraic over F. Then

$$E = F(\alpha_1, \ldots, \alpha_n) \supset F(\alpha_1, \ldots, \alpha_{n-1}) \supset \cdots \supset F(\alpha_1) \supset F,$$

where each field $F(\alpha_1, \ldots, \alpha_i)$ is algebraic over $F(\alpha_1, \ldots, \alpha_{i-1})$.

$(3) \Rightarrow (1)$. Let

$$E = F(\alpha_1, \ldots, \alpha_n) \supset F(\alpha_1, \ldots, \alpha_{n-1}) \supset \cdots \supset F(\alpha_1) \supset F,$$

where each field $F(\alpha_1, \ldots, \alpha_i)$ is algebraic over $F(\alpha_1, \ldots, \alpha_{i-1})$. Since

$$F(\alpha_1, \ldots, \alpha_i) = F(\alpha_1, \ldots, \alpha_{i-1})(\alpha_i)$$

is simple extension and α_i is algebraic over $F(\alpha_1, \ldots, \alpha_{i-1})$, it follows that

$$[F(\alpha_1, \ldots, \alpha_i) : F(\alpha_1, \ldots, \alpha_{i-1})]$$

is finite for each i. Therefore, $[E : F]$ is finite. $\qquad \square$

Algebraic Closure

Given a field F, the question arises as to whether or not we can find a field E such that every polynomial $p(x)$ has a root in E. This leads us to the following theorem.

Theorem 21.23. *Let E be an extension field of F. The set of elements in E that are algebraic over F form a field.*

PROOF. Let $\alpha, \beta \in E$ be algebraic over F. Then $F(\alpha, \beta)$ is a finite extension of F. Since every element of $F(\alpha, \beta)$ is algebraic over F, $\alpha \pm \beta$, $\alpha\beta$, and α/β ($\beta \neq 0$) are all algebraic over F. Consequently, the set of elements in E that are algebraic over F form a field. $\quad \square$

Corollary 21.24. *The set of all algebraic numbers forms a field; that is, the set of all complex numbers that are algebraic over $\mathbb{Q}$ makes up a field.*

Let E be a field extension of a field F. We define the **algebraic closure** of a field F in E to be the field consisting of all elements in E that are algebraic over F. A field F is **algebraically closed** if every nonconstant polynomial in $F[x]$ has a root in F.

Theorem 21.25. *A field F is algebraically closed if and only if every nonconstant polynomial in $F[x]$ factors into linear factors over $F[x]$.*

PROOF.   Let F be an algebraically closed field. If $p(x) \in F[x]$ is a nonconstant polynomial, then $p(x)$ has a zero in F, say α. Therefore, $x - \alpha$ must be a factor of $p(x)$ and so $p(x) = (x - \alpha)q_1(x)$, where $\deg q_1(x) = \deg p(x) - 1$. Continue this process with $q_1(x)$ to find a factorization

$$p(x) = (x - \alpha)(x - \beta)q_2(x),$$

where $\deg q_2(x) = \deg p(x) - 2$. The process must eventually stop since the degree of $p(x)$ is finite.

Conversely, suppose that every nonconstant polynomial $p(x)$ in $F[x]$ factors into linear factors. Let $ax - b$ be such a factor. Then $p(b/a) = 0$. Consequently, F is algebraically closed.   □

Corollary 21.26. *An algebraically closed field F has no proper algebraic extension E.*

PROOF.   Let E be an algebraic extension of F; then $F \subset E$. For $\alpha \in E$, the minimal polynomial of α is $x - \alpha$. Therefore, $\alpha \in F$ and $F = E$.   □

Theorem 21.27. *Every field F has a unique algebraic closure.*

It is a nontrivial fact that every field has a unique algebraic closure. The proof is not extremely difficult, but requires some rather sophisticated set theory. We refer the reader to [3], [4], or [8] for a proof of this result.

We now state the Fundamental Theorem of Algebra, first proven by Gauss at the age of 22 in his doctoral thesis. This theorem states that every polynomial with coefficients in the complex numbers has a root in the complex numbers. The proof of this theorem will be given in Chapter 23.

Theorem 21.28 Fundamental Theorem of Algebra. *The field of complex numbers is algebraically closed.*

21.2 Splitting Fields

Let F be a field and $p(x)$ be a nonconstant polynomial in $F[x]$. We already know that we can find a field extension of F that contains a root of $p(x)$. However, we would like to know whether an extension E of F containing all of the roots of $p(x)$ exists. In other words, can we find a field extension of F such that $p(x)$ factors into a product of linear polynomials? What is the "smallest" extension containing all the roots of $p(x)$?

Let F be a field and $p(x) = a_0 + a_1 x + \cdots + a_n x^n$ be a nonconstant polynomial in $F[x]$. An extension field E of F is a **splitting field** of $p(x)$ if there exist elements $\alpha_1, \ldots, \alpha_n$ in E such that $E = F(\alpha_1, \ldots, \alpha_n)$ and

$$p(x) = (x - \alpha_1)(x - \alpha_2) \cdots (x - \alpha_n).$$

A polynomial $p(x) \in F[x]$ **splits** in E if it is the product of linear factors in $E[x]$.

Example 21.29. Let $p(x) = x^4 + 2x^2 - 8$ be in $\mathbb{Q}[x]$. Then $p(x)$ has irreducible factors $x^2 - 2$ and $x^2 + 4$. Therefore, the field $\mathbb{Q}(\sqrt{2}, i)$ is a splitting field for $p(x)$.

Example 21.30. Let $p(x) = x^3 - 3$ be in $\mathbb{Q}[x]$. Then $p(x)$ has a root in the field $\mathbb{Q}(\sqrt[3]{3})$. However, this field is not a splitting field for $p(x)$ since the complex cube roots of 3,

$$\frac{-\sqrt[3]{3} \pm (\sqrt[6]{3})^5 i}{2},$$

are not in $\mathbb{Q}(\sqrt[3]{3})$.

Theorem 21.31. *Let $p(x) \in F[x]$ be a nonconstant polynomial. Then there exists a splitting field E for $p(x)$.*

PROOF.   We will use mathematical induction on the degree of $p(x)$. If $\deg p(x) = 1$, then $p(x)$ is a linear polynomial and $E = F$. Assume that the theorem is true for all polynomials of degree k with $1 \leq k < n$ and let $\deg p(x) = n$. We can assume that $p(x)$ is irreducible; otherwise, by our induction hypothesis, we are done. By Theorem 21.5, there exists a field K such that $p(x)$ has a zero α_1 in K. Hence, $p(x) = (x - \alpha_1)q(x)$, where $q(x) \in K[x]$. Since $\deg q(x) = n - 1$, there exists a splitting field $E \supset K$ of $q(x)$ that contains the zeros $\alpha_2, \ldots, \alpha_n$ of $p(x)$ by our induction hypothesis. Consequently,

$$E = K(\alpha_2, \ldots, \alpha_n) = F(\alpha_1, \ldots, \alpha_n)$$

is a splitting field of $p(x)$. $\square$

The question of uniqueness now arises for splitting fields. This question is answered in the affirmative. Given two splitting fields K and L of a polynomial $p(x) \in F[x]$, there exists a field isomorphism $\phi : K \to L$ that preserves F. In order to prove this result, we must first prove a lemma.

Lemma 21.32. *Let $\phi : E \to F$ be an isomorphism of fields. Let K be an extension field of E and $\alpha \in K$ be algebraic over E with minimal polynomial $p(x)$. Suppose that L is an extension field of F such that β is root of the polynomial in $F[x]$ obtained from $p(x)$ under the image of ϕ. Then ϕ extends to a unique isomorphism $\overline{\phi} : E(\alpha) \to F(\beta)$ such that $\overline{\phi}(\alpha) = \beta$ and $\overline{\phi}$ agrees with ϕ on E.*

PROOF.   If $p(x)$ has degree n, then by Theorem 21.13 we can write any element in $E(\alpha)$ as a linear combination of $1, \alpha, \ldots, \alpha^{n-1}$. Therefore, the isomorphism that we are seeking must be

$$\overline{\phi}(a_0 + a_1\alpha + \cdots + a_{n-1}\alpha^{n-1}) = \phi(a_0) + \phi(a_1)\beta + \cdots + \phi(a_{n-1})\beta^{n-1},$$

where

$$a_0 + a_1\alpha + \cdots + a_{n-1}\alpha^{n-1}$$

is an element in $E(\alpha)$. The fact that $\overline{\phi}$ is an isomorphism could be checked by direct computation; however, it is easier to observe that $\overline{\phi}$ is a composition of maps that we already know to be isomorphisms.

We can extend ϕ to be an isomorphism from $E[x]$ to $F[x]$, which we will also denote by ϕ, by letting

$$\phi(a_0 + a_1 x + \cdots + a_n x^n) = \phi(a_0) + \phi(a_1)x + \cdots + \phi(a_n)x^n.$$

This extension agrees with the original isomorphism $\phi : E \to F$, since constant polynomials get mapped to constant polynomials. By assumption, $\phi(p(x)) = q(x)$; hence, ϕ maps $\langle p(x) \rangle$ onto $\langle q(x) \rangle$. Consequently, we have an isomorphism $\psi : E[x]/\langle p(x) \rangle \to F[x]/\langle q(x) \rangle$. By Proposition 21.12, we have isomorphisms $\sigma : E[x]/\langle p(x) \rangle \to E(\alpha)$ and $\tau : F[x]/\langle q(x) \rangle \to F(\beta)$, defined by evaluation at α and β, respectively. Therefore, $\overline{\phi} = \tau\psi\sigma^{-1}$ is the required isomorphism.

$$E[x]/\langle p(x)\rangle \xrightarrow{\ \psi\ } F[x]/\langle q(x)\rangle$$

$$\downarrow \sigma \qquad\qquad\qquad \downarrow \tau$$

$$E(\alpha) \xrightarrow{\ \overline{\phi}\ } F(\beta)$$

$$\downarrow \qquad\qquad\qquad \downarrow$$

$$E \xrightarrow{\ \phi\ } F$$

We leave the proof of uniqueness as a exercise. $\qquad\qquad\qquad\qquad\qquad\square$

Theorem 21.33. *Let $\phi : E \to F$ be an isomorphism of fields and let $p(x)$ be a nonconstant polynomial in $E[x]$ and $q(x)$ the corresponding polynomial in $F[x]$ under the isomorphism. If K is a splitting field of $p(x)$ and L is a splitting field of $q(x)$, then ϕ extends to an isomorphism $\psi : K \to L$.*

PROOF. We will use mathematical induction on the degree of $p(x)$. We can assume that $p(x)$ is irreducible over E. Therefore, $q(x)$ is also irreducible over F. If $\deg p(x) = 1$, then by the definition of a splitting field, $K = E$ and $L = F$ and there is nothing to prove.

Assume that the theorem holds for all polynomials of degree less than n. Since K is a splitting field of $p(x)$, all of the roots of $p(x)$ are in K. Choose one of these roots, say α, such that $E \subset E(\alpha) \subset K$. Similarly, we can find a root β of $q(x)$ in L such that $F \subset F(\beta) \subset L$. By Lemma 21.32, there exists an isomorphism $\overline{\phi} : E(\alpha) \to F(\beta)$ such that $\overline{\phi}(\alpha) = \beta$ and $\overline{\phi}$ agrees with ϕ on E.

$$K \xrightarrow{\ \psi\ } L$$

$$\downarrow \qquad\qquad\qquad \downarrow$$

$$E(\alpha) \xrightarrow{\ \overline{\phi}\ } F(\beta)$$

$$\downarrow \qquad\qquad\qquad \downarrow$$

$$E \xrightarrow{\ \phi\ } F$$

Now write $p(x) = (x - \alpha)f(x)$ and $q(x) = (x - \beta)g(x)$, where the degrees of $f(x)$ and $g(x)$ are less than the degrees of $p(x)$ and $q(x)$, respectively. The field extension K is a

splitting field for $f(x)$ over $E(\alpha)$, and L is a splitting field for $g(x)$ over $F(\beta)$. By our induction hypothesis there exists an isomorphism $\psi : K \to L$ such that ψ agrees with $\overline{\phi}$ on $E(\alpha)$. Hence, there exists an isomorphism $\psi : K \to L$ such that ψ agrees with ϕ on E. $\square$

Corollary 21.34. *Let $p(x)$ be a polynomial in $F[x]$. Then there exists a splitting field K of $p(x)$ that is unique up to isomorphism.*

21.3 Geometric Constructions

In ancient Greece, three classic problems were posed. These problems are geometric in nature and involve straightedge-and-compass constructions from what is now high school geometry; that is, we are allowed to use only a straightedge and compass to solve them. The problems can be stated as follows.

1. Given an arbitrary angle, can one trisect the angle into three equal subangles using only a straightedge and compass?

2. Given an arbitrary circle, can one construct a square with the same area using only a straightedge and compass?

3. Given a cube, can one construct the edge of another cube having twice the volume of the original? Again, we are only allowed to use a straightedge and compass to do the construction.

After puzzling mathematicians for over two thousand years, each of these constructions was finally shown to be impossible. We will use the theory of fields to provide a proof that the solutions do not exist. It is quite remarkable that the long-sought solution to each of these three geometric problems came from abstract algebra.

First, let us determine more specifically what we mean by a straightedge and compass, and also examine the nature of these problems in a bit more depth. To begin with, *a straightedge is not a ruler*. We cannot measure arbitrary lengths with a straightedge. It is merely a tool for drawing a line through two points. The statement that the trisection of an arbitrary angle is impossible means that there is at least one angle that is impossible to trisect with a straightedge-and-compass construction. Certainly it is possible to trisect an angle in special cases. We can construct a $30°$ angle; hence, it is possible to trisect a $90°$ angle. However, we will show that it is impossible to construct a $20°$ angle. Therefore, we cannot trisect a $60°$ angle.

Constructible Numbers

A real number α is **constructible** if we can construct a line segment of length $|\alpha|$ in a finite number of steps from a segment of unit length by using a straightedge and compass.

Theorem 21.35. *The set of all constructible real numbers forms a subfield F of the field of real numbers.*

PROOF. Let α and β be constructible numbers. We must show that $\alpha + \beta$, $\alpha - \beta$, $\alpha\beta$, and α/β ($\beta \neq 0$) are also constructible numbers. We can assume that both α and β are positive with $\alpha > \beta$. It is quite obvious how to construct $\alpha + \beta$ and $\alpha - \beta$. To find a line segment with length $\alpha\beta$, we assume that $\beta > 1$ and construct the triangle in Figure 21.36 such that triangles $\triangle ABC$ and $\triangle ADE$ are similar. Since $\alpha/1 = x/\beta$, the line segment x has length $\alpha\beta$. A similar construction can be made if $\beta < 1$. We will leave it as an exercise to show that the same triangle can be used to construct α/β for $\beta \neq 0$. $\square$
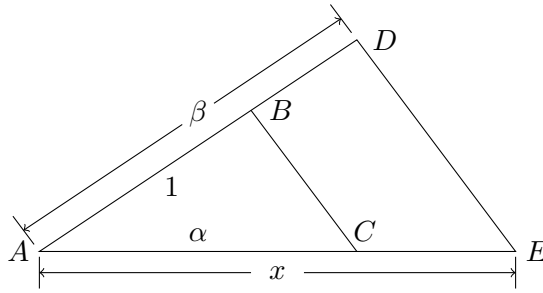
Figure 21.36: Construction of products

Lemma 21.37. *If α is a constructible number, then $\sqrt{\alpha}$ is a constructible number.*

PROOF.   In Figure 21.38 the triangles $\triangle ABD$, $\triangle BCD$, and $\triangle ABC$ are similar; hence, $1/x = x/\alpha$, or $x^2 = \alpha$. $\qquad\square$
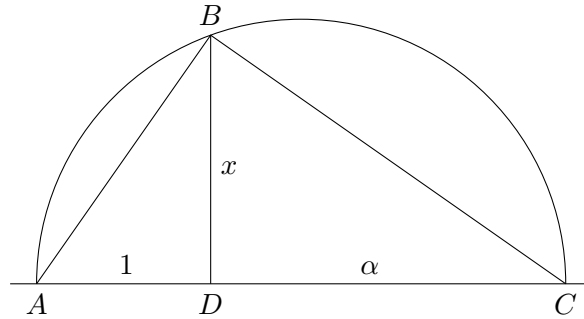
Figure 21.38: Construction of roots

By Theorem 21.35, we can locate in the plane any point $P = (p, q)$ that has rational coordinates p and q. We need to know what other points can be constructed with a compass and straightedge from points with rational coordinates.

Lemma 21.39. *Let F be a subfield of $\mathbb{R}$.*

1. *If a line contains two points in F, then it has the equation $ax + by + c = 0$, where a, b, and c are in F.*

2. *If a circle has a center at a point with coordinates in F and a radius that is also in F, then it has the equation $x^2 + y^2 + dx + ey + f = 0$, where d, e, and f are in F.*

PROOF.   Let (x_1, y_1) and (x_2, y_2) be points on a line whose coordinates are in F. If $x_1 = x_2$, then the equation of the line through the two points is $x - x_1 = 0$, which has the form $ax + by + c = 0$. If $x_1 \neq x_2$, then the equation of the line through the two points is given by

$$y - y_1 = \left(\frac{y_2 - y_1}{x_2 - x_1} \right) (x - x_1),$$

which can also be put into the proper form.

To prove the second part of the lemma, suppose that (x_1, y_1) is the center of a circle of radius r. Then the circle has the equation

$$(x - x_1)^2 + (y - y_1)^2 - r^2 = 0.$$

This equation can easily be put into the appropriate form. $\qquad\square$

Starting with a field of constructible numbers F, we have three possible ways of constructing additional points in $\mathbb{R}$ with a compass and straightedge.

1. To find possible new points in $\mathbb{R}$, we can take the intersection of two lines, each of which passes through two known points with coordinates in F.

2. The intersection of a line that passes through two points that have coordinates in F and a circle whose center has coordinates in F with radius of a length in F will give new points in $\mathbb{R}$.

3. We can obtain new points in $\mathbb{R}$ by intersecting two circles whose centers have coordinates in F and whose radii are of lengths in F.

The first case gives no new points in $\mathbb{R}$, since the solution of two equations of the form $ax + by + c = 0$ having coefficients in F will always be in F. The third case can be reduced to the second case. Let

$$x^2 + y^2 + d_1 x + e_1 y + f_1 = 0$$
$$x^2 + y^2 + d_2 x + e_2 y + f_2 = 0$$

be the equations of two circles, where d_i, e_i, and f_i are in F for $i = 1, 2$. These circles have the same intersection as the circle

$$x^2 + y^2 + d_1 x + e_1 x + f_1 = 0$$

and the line

$$(d_1 - d_2)x + b(e_2 - e_1)y + (f_2 - f_1) = 0.$$

The last equation is that of the chord passing through the intersection points of the two circles. Hence, the intersection of two circles can be reduced to the case of an intersection of a line with a circle.

Considering the case of the intersection of a line and a circle, we must determine the nature of the solutions of the equations

$$ax + by + c = 0$$
$$x^2 + y^2 + dx + ey + f = 0.$$

If we eliminate y from these equations, we obtain an equation of the form $Ax^2 + Bx + C = 0$, where A, B, and C are in F. The x coordinate of the intersection points is given by

$$x = \frac{-B \pm \sqrt{B^2 - 4AC}}{2A}$$

and is in $F(\sqrt{\alpha})$, where $\alpha = B^2 - 4AC > 0$. We have proven the following lemma.

Lemma 21.40. *Let F be a field of constructible numbers. Then the points determined by the intersections of lines and circles in F lie in the field $F(\sqrt{\alpha})$ for some α in F.*

Theorem 21.41. *A real number α is a constructible number if and only if there exists a sequence of fields*

$$\mathbb{Q} = F_0 \subset F_1 \subset \cdots \subset F_k$$

such that $F_i = F_{i-1}(\sqrt{\alpha_i})$ with $\alpha_i \in F_i$ and $\alpha \in F_k$. In particular, there exists an integer $k > 0$ such that $[\mathbb{Q}(\alpha) : \mathbb{Q}] = 2^k$.

PROOF.   The existence of the F_i's and the α_i's is a direct consequence of Lemma 21.40 and of the fact that

$$[F_k : \mathbb{Q}] = [F_k : F_{k-1}][F_{k-1} : F_{k-2}] \cdots [F_1 : \mathbb{Q}] = 2^k. \qquad \square$$

Corollary 21.42. *The field of all constructible numbers is an algebraic extension of $\mathbb{Q}$.*

As we can see by the field of constructible numbers, not every algebraic extension of a field is a finite extension.

Doubling the Cube and Squaring the Circle

We are now ready to investigate the classical problems of doubling the cube and squaring the circle. We can use the field of constructible numbers to show exactly when a particular geometric construction can be accomplished.

Doubling the cube is impossible   Given the edge of the cube, it is impossible to construct with a straightedge and compass the edge of the cube that has twice the volume of the original cube. Let the original cube have an edge of length 1 and, therefore, a volume of 1. If we could construct a cube having a volume of 2, then this new cube would have an edge of length $\sqrt[3]{2}$. However, $\sqrt[3]{2}$ is a zero of the irreducible polynomial $x^3 - 2$ over $\mathbb{Q}$; hence,

$$[\mathbb{Q}(\sqrt[3]{2}) : \mathbb{Q}] = 3$$

This is impossible, since 3 is not a power of 2.

Squaring the circle   Suppose that we have a circle of radius 1. The area of the circle is π; therefore, we must be able to construct a square with side $\sqrt{\pi}$. This is impossible since π and consequently $\sqrt{\pi}$ are both transcendental. Therefore, using a straightedge and compass, it is not possible to construct a square with the same area as the circle.

Trisecting an Angle

Trisecting an arbitrary angle is impossible. We will show that it is impossible to construct a 20° angle. Consequently, a 60° angle cannot be trisected. We first need to calculate the triple-angle formula for the cosine:

$$\begin{aligned}
\cos 3\theta &= \cos(2\theta + \theta) \\
&= \cos 2\theta \cos \theta - \sin 2\theta \sin \theta \\
&= (2\cos^2 \theta - 1)\cos \theta - 2\sin^2 \theta \cos \theta \\
&= (2\cos^2 \theta - 1)\cos \theta - 2(1 - \cos^2 \theta)\cos \theta \\
&= 4\cos^3 \theta - 3\cos \theta.
\end{aligned}$$

The angle θ can be constructed if and only if $\alpha = \cos \theta$ is constructible. Let $\theta = 20°$. Then $\cos 3\theta = \cos 60° = 1/2$. By the triple-angle formula for the cosine,

$$4\alpha^3 - 3\alpha = \frac{1}{2}.$$

Therefore, α is a zero of $8x^3 - 6x - 1$. This polynomial has no factors in $\mathbb{Z}[x]$, and hence is irreducible over $\mathbb{Q}[x]$. Thus, $[\mathbb{Q}(\alpha) : \mathbb{Q}] = 3$. Consequently, α cannot be a constructible number.

Algebraic number theory uses the tools of algebra to solve problems in number theory. Modern algebraic number theory began with Pierre de Fermat (1601–1665). Certainly we can find many positive integers that satisfy the equation $x^2 + y^2 = z^2$; Fermat conjectured that the equation $x^n + y^n = z^n$ has no positive integer solutions for $n \geq 3$. He stated in the margin of his copy of the Latin translation of Diophantus' *Arithmetica* that he had found a marvelous proof of this theorem, but that the margin of the book was too narrow to contain it. Building on work of other mathematicians, it was Andrew Wiles who finally succeeded in proving Fermat's Last Theorem in the 1990s. Wiles's achievement was reported on the front page of the *New York Times*.

Attempts to prove Fermat's Last Theorem have led to important contributions to algebraic number theory by such notable mathematicians as Leonhard Euler (1707–1783). Significant advances in the understanding of Fermat's Last Theorem were made by Ernst Kummer (1810–1893). Kummer's student, Leopold Kronecker (1823–1891), became one of the leading algebraists of the nineteenth century. Kronecker's theory of ideals and his study of algebraic number theory added much to the understanding of fields.

David Hilbert (1862–1943) and Hermann Minkowski (1864–1909) were among the mathematicians who led the way in this subject at the beginning of the twentieth century. Hilbert and Minkowski were both mathematicians at Göttingen University in Germany. Göttingen was truly one the most important centers of mathematical research during the last two centuries. The large number of exceptional mathematicians who studied there included Gauss, Dirichlet, Riemann, Dedekind, Noether, and Weyl.

André Weil answered questions in number theory using algebraic geometry, a field of mathematics that studies geometry by studying commutative rings. From about 1955 to 1970, Alexander Grothendieck dominated the field of algebraic geometry. Pierre Deligne, a student of Grothendieck, solved several of Weil's number-theoretic conjectures. One of the most recent contributions to algebra and number theory is Gerd Falting's proof of the Mordell-Weil conjecture. This conjecture of Mordell and Weil essentially says that certain polynomials $p(x, y)$ in $\mathbb{Z}[x, y]$ have only a finite number of integral solutions.

21.4 Exercises

1. Show that each of the following numbers is algebraic over $\mathbb{Q}$ by finding the minimal polynomial of the number over $\mathbb{Q}$.

(a) $\sqrt{1/3 + \sqrt{7}}$

(b) $\sqrt{3} + \sqrt[3]{5}$

(c) $\sqrt{3} + \sqrt{2}\, i$

(d) $\cos\theta + i\sin\theta$ for $\theta = 2\pi/n$ with $n \in \mathbb{N}$

(e) $\sqrt{\sqrt[3]{2} - i}$

2. Find a basis for each of the following field extensions. What is the degree of each extension?

(a) $\mathbb{Q}(\sqrt{3}, \sqrt{6})$ over $\mathbb{Q}$

(b) $\mathbb{Q}(\sqrt[3]{2}, \sqrt[3]{3})$ over $\mathbb{Q}$

(c) $\mathbb{Q}(\sqrt{2}, i)$ over $\mathbb{Q}$

(d) $\mathbb{Q}(\sqrt{3}, \sqrt{5}, \sqrt{7})$ over $\mathbb{Q}$

(e) $\mathbb{Q}(\sqrt{2}, \sqrt[3]{2})$ over $\mathbb{Q}$

(f) $\mathbb{Q}(\sqrt{8})$ over $\mathbb{Q}(\sqrt{2})$

(g) $\mathbb{Q}(i, \sqrt{2}+i, \sqrt{3}+i)$ over $\mathbb{Q}$

(h) $\mathbb{Q}(\sqrt{2}+\sqrt{5})$ over $\mathbb{Q}(\sqrt{5})$

(i) $\mathbb{Q}(\sqrt{2}, \sqrt{6}+\sqrt{10})$ over $\mathbb{Q}(\sqrt{3}+\sqrt{5})$

3.   Find the splitting field for each of the following polynomials.

(a) $x^4 - 10x^2 + 21$ over $\mathbb{Q}$

(b) $x^4 + 1$ over $\mathbb{Q}$

(c) $x^3 + 2x + 2$ over $\mathbb{Z}_3$

(d) $x^3 - 3$ over $\mathbb{Q}$

4.   Consider the field extension $\mathbb{Q}(\sqrt[4]{3}, i)$ over $\mathbb{Q}$.

(a) Find a basis for the field extension $\mathbb{Q}(\sqrt[4]{3}, i)$ over $\mathbb{Q}$. Conclude that $[\mathbb{Q}(\sqrt[4]{3}, i) : \mathbb{Q}] = 8$.

(b) Find all subfields F of $\mathbb{Q}(\sqrt[4]{3}, i)$ such that $[F : \mathbb{Q}] = 2$.

(c) Find all subfields F of $\mathbb{Q}(\sqrt[4]{3}, i)$ such that $[F : \mathbb{Q}] = 4$.

5.   Show that $\mathbb{Z}_2[x]/\langle x^3 + x + 1 \rangle$ is a field with eight elements. Construct a multiplication table for the multiplicative group of the field.

6.   Show that the regular 9-gon is not constructible with a straightedge and compass, but that the regular 20-gon is constructible.

7.   Prove that the cosine of one degree $(\cos 1°)$ is algebraic over $\mathbb{Q}$ but not constructible.

8.   Can a cube be constructed with three times the volume of a given cube?

9.   Prove that $\mathbb{Q}(\sqrt{3}, \sqrt[4]{3}, \sqrt[8]{3}, \ldots)$ is an algebraic extension of $\mathbb{Q}$ but not a finite extension.

10.   Prove or disprove: π is algebraic over $\mathbb{Q}(\pi^3)$.

11.   Let $p(x)$ be a nonconstant polynomial of degree n in $F[x]$. Prove that there exists a splitting field E for $p(x)$ such that $[E : F] \leq n!$.

12.   Prove or disprove: $\mathbb{Q}(\sqrt{2}) \cong \mathbb{Q}(\sqrt{3})$.

13.   Prove that the fields $\mathbb{Q}(\sqrt[4]{3})$ and $\mathbb{Q}(\sqrt[4]{3}\,i)$ are isomorphic but not equal.

14.   Let K be an algebraic extension of E, and E an algebraic extension of F. Prove that K is algebraic over F. [*Caution*: Do not assume that the extensions are finite.]

15.   Prove or disprove: $\mathbb{Z}[x]/\langle x^3 - 2 \rangle$ is a field.

16.   Let F be a field of characteristic p. Prove that $p(x) = x^p - a$ either is irreducible over F or splits in F.

17.   Let E be the algebraic closure of a field F. Prove that every polynomial $p(x)$ in $F[x]$ splits in E.

18.   If every irreducible polynomial $p(x)$ in $F[x]$ is linear, show that F is an algebraically closed field.

19.   Prove that if α and β are constructible numbers such that $\beta \neq 0$, then so is α/β.

20.   Show that the set of all elements in $\mathbb{R}$ that are algebraic over $\mathbb{Q}$ form a field extension of $\mathbb{Q}$ that is not finite.

21.   Let E be an algebraic extension of a field F, and let σ be an automorphism of E leaving F fixed. Let $\alpha \in E$. Show that σ induces a permutation of the set of all zeros of the minimal polynomial of α that are in E.

22.   Show that $\mathbb{Q}(\sqrt{3}, \sqrt{7}) = \mathbb{Q}(\sqrt{3} + \sqrt{7})$. Extend your proof to show that $\mathbb{Q}(\sqrt{a}, \sqrt{b}) = \mathbb{Q}(\sqrt{a} + \sqrt{b})$, where $\gcd(a, b) = 1$.

23.   Let E be a finite extension of a field F. If $[E : F] = 2$, show that E is a splitting field of F for some polynomial $f(x) \in F[x]$.

24.   Prove or disprove: Given a polynomial $p(x)$ in $\mathbb{Z}_6[x]$, it is possible to construct a ring R such that $p(x)$ has a root in R.

25.   Let E be a field extension of F and $\alpha \in E$. Determine $[F(\alpha) : F(\alpha^3)]$.

26.   Let α, β be transcendental over $\mathbb{Q}$. Prove that either $\alpha\beta$ or $\alpha+\beta$ is also transcendental.

27.   Let E be an extension field of F and $\alpha \in E$ be transcendental over F. Prove that every element in $F(\alpha)$ that is not in F is also transcendental over F.

28.   Let α be a root of an irreducible monic polynomial $p(x) \in F[x]$, with $\deg p = n$. Prove that $[F(\alpha) : F] = n$.

References and Suggested Readings

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21.5   Sage

In Sage, and other places, an extension of the rationals is called a "number field." They are one of Sage's most mature features.

Number Fields

There are several ways to create a number field. We are familiar with the syntax where we adjoin an irrational number that we can write with traditional combinations of arithmetic and roots.

```
M.<a> = QQ[sqrt(2)+sqrt(3)]; M
```

```
Number Field in a with defining polynomial x^4 - 10*x^2 + 1
```

We can also specify the element we want to adjoin as the root of a monic irreducible polynomial. One approach is to construct the polynomial ring first so that the polynomial has the location of its coefficients specified properly.

```
F.<y> = QQ[]
p = y^3 - 1/4*y^2 - 1/16*y + 1/4
p.is_irreducible()
```

```
True
```

```
N.<b> = NumberField(p, 'b'); N
```

```
Number Field in b with
defining polynomial y^3 - 1/4*y^2 - 1/16*y + 1/4
```

Rather than building the whole polynomial ring, we can simply introduce a variable as the generator of a polynomial ring and then create polynomials from this variable. This spares us naming the polynomial ring. Notice in the example that both instances of `z` are necessary.

```
z = polygen(QQ, 'z')
q = z^3 - 1/4*z^2 - 1/16*z + 1/4
q.parent()
```

```
Univariate Polynomial Ring in z over Rational Field
```

```
P.<c> = NumberField(q, 'c'); P
```

```
Number Field in c with
defining polynomial z^3 - 1/4*z^2 - 1/16*z + 1/4
```

We can recover the polynomial used to create a number field, even if we constructed it by giving an expression for an irrational element. In this case, the polynomial is the minimal polynomial of the element.

```
M.polynomial()
```

```
x^4 - 10*x^2 + 1
```

```
N.polynomial()
```

```
y^3 - 1/4*y^2 - 1/16*y + 1/4
```

For any element of a number field, Sage will obligingly compute its minimal polynomial.

```
element = -b^2 + 1/3*b + 4
element.parent()
```

```
Number Field in b with
defining polynomial y^3 - 1/4*y^2 - 1/16*y + 1/4
```

```
r = element.minpoly('t'); r
```

```
t^3 - 571/48*t^2 + 108389/2304*t - 13345/216
```

```
r.parent()
```

```
Univariate Polynomial Ring in t over Rational Field
```

```
r.subs(t=element)
```

```
0
```

Substituting `element` back into the alleged minimal polynomial and getting back zero is not convincing evidence that it is the *minimal* polynomial, but it is heartening.

Relative and Absolute Number Fields

With Sage we can adjoin several elements at once and we can build nested towers of number fields. Sage uses the term "absolute" to refer to a number field viewed as an extension of the rationals themselves, and the term "relative" to refer to a number field constructed, or viewed, as an extension of another (nontrivial) number field.

```
A.<a,b> = QQ[sqrt(2), sqrt(3)]
A
```

```
Number Field in sqrt2 with defining polynomial x^2 - 2 over
its base field
```

```
B = A.base_field(); B
```

```
Number Field in sqrt3 with defining polynomial x^2 - 3
```

```
A.is_relative()
```

```
True
```

```
B.is_relative()
```

```
False
```

The number field `A` has been constructed mathematically as what we would write as $\mathbb{Q} \subset \mathbb{Q}[\sqrt{3}] \subset \mathbb{Q}[\sqrt{3}, \sqrt{2}]$. Notice the slight difference in ordering of the elements we are adjoining, and notice how the number fields use slightly fancier internal names (`sqrt2`, `sqrt3`) for the new elements.

We can "flatten" a relative field to view it as an absolute field, which may have been our intention from the start. Here we create a new number field from `A` that makes it a pure absolute number field.

```
C.<c> = A.absolute_field()
C
```

```
Number Field in c with defining polynomial x^4 - 10*x^2 + 1
```

Once we construct an absolute number field this way, we can recover isomorphisms to and from the absolute field. Recall that our tower was built with generators `a` and `b`, while the flattened tower is generated by `c`. The `.structure()` method returns a pair of functions, with the absolute number field as the domain and codomain (in that order).

```
fromC, toC = C.structure()
fromC(c)
```

```
sqrt2 - sqrt3
```

```
toC(a)
```

```
1/2*c^3 - 9/2*c
```

```
toC(b)
```

```
1/2*c^3 - 11/2*c
```

This tells us that the single generator of the flattened tower, `c`, is equal to $\sqrt{2} - \sqrt{3}$, and further, each of $\sqrt{2}$ and $\sqrt{3}$ can be expressed as polynomial functions of `c`. With these connections, you might want to compute the final two expressions in `c` by hand, and appreciate the work Sage does to determine these for us. This computation is an example of the conclusion of the upcoming Theorem 23.12.

Many number field methods have both relative and absolute versions, and we will also find it more convenient to work in a tower or a flattened version, thus the isomorphisms between the two can be invaluable for translating both questions and answers.

As a vector space over $\mathbb{Q}$, or over another number field, number fields that are finite extensions have a dimension, called the degree. These are easy to get from Sage, though for a relative field, we need to be more precise about which degree we desire.

```
B.degree()
```

```
2
```

```
A.absolute_degree()
```

```
4
```

```
A.relative_degree()
```

```
2
```

Splitting Fields

Here is a concrete example of how to use Sage to construct a splitting field of a polynomial. Consider $p(x) = x^4 + x^2 - 1$. We first build a number field with a single root, and then factor the polynomial over this new, larger, field.

```
x = polygen(QQ, 'x')
p = x^4 + x^2 - 1
p.parent()
```

```
Univariate Polynomial Ring in x over Rational Field
```

```
p.is_irreducible()
```

```
True
```

```
M.<a> = NumberField(p, 'a')
y = polygen(M, 'y')
p = p.subs(x = y)
p
```

```
y^4 + y^2 - 1
```

```
p.parent()
```

```
Univariate Polynomial Ring in y over Number Field in a with
defining polynomial x^4 + x^2 - 1
```

```
p.factor()
```

```
(y - a) * (y + a) * (y^2 + a^2 + 1)
```

```
a^2 + 1 in QQ
```

```
False
```

So our polynomial factors partially into two linear factors and a quadratic factor. But notice that the quadratic factor has a coefficient that is irrational, $a^2 + 1$, so the quadratic factor properly belongs in the polynomial ring over M and not over QQ.

We build an extension containing a root of the quadratic factor, called q here. Then, rather than using the polygen() function, we build an entire polynomial ring R over N with the indeterminate z. The reason for doing this is we can illustrate how we "upgrade" the polynomial p with the syntax R(p) to go from having coefficients in M to having coefficients in N.

```
q = y^2 + a^2 + 1
N.<b> = NumberField(q, 'b')
R.<z> = N[]
s = R(p)
s
```

```
z^4 + z^2 - 1
```

```
s.parent()
```

```
Univariate Polynomial Ring in z over Number Field in b with
defining polynomial y^2 + a^2 + 1 over its base field
```

```
s.factor()
```

```
(z + b) * (z + a) * (z - a) * (z - b)
```

```
a in N, b in N
```

```
(True, True)
```

So we have a field, N, where our polynomial factors into linear factors with coefficients from the field. We can get another factorization by converting N to an absolute number field and factoring there. We need to recreate the polynomial over N, since a substitution will carry coefficients from the wrong ring.

```
P.<c> = N.absolute_field()
w = polygen(P, 'w')
p = w^4 + w^2- 1
p.factor()
```

```
(w - 7/18966*c^7 + 110/9483*c^5 +  923/9483*c^3 +  3001/6322*c) *
(w - 7/37932*c^7 +  55/9483*c^5 + 923/18966*c^3 - 3321/12644*c) *
(w + 7/37932*c^7 -  55/9483*c^5 - 923/18966*c^3 + 3321/12644*c) *
(w + 7/18966*c^7 - 110/9483*c^5 -  923/9483*c^3 -  3001/6322*c)
```

This is an interesting alternative, in that the roots of the polynomial are expressions in terms of the *single* generator `c`. Since the roots involve a seventh power of `c`, we might suspect (but not be certain) that the minimal polynomial of `c` has degree 8 and that `P` is a degree 8 extension of the rationals. Indeed `P` (or `N`) is a splitting field for $p(x) = x^4 + x^2 - 1$. The roots are not really as bad as they appear — lets convert them back to the relative number field.

First we want to rewrite a single factor (the first) in the form $(w - r)$ to identify the root with the correct signs.

```
(w -   7/18966*c^7 + 110/9483*c^5 + 923/9483*c^3 + 3001/6322*c) =
(w - (7/18966*c^7 - 110/9483*c^5 - 923/9483*c^3 - 3001/6322*c))
```

With the conversion isomorphisms, we can recognize the roots for what they are.

```
fromP, toP = P.structure()
fromP(7/18966*c^7 - 110/9483*c^5 - 923/9483*c^3 - 3001/6322*c)
```

```
-b
```

So the rather complicated expression in `c` is just the negative of the root we adjoined in the second step of constructing the tower of number fields. It would be a good exercise to see what happens to the other three roots (being careful to get the signs right on each root).

This is a good opportunity to illustrate Theorem 21.17.

```
M.degree()
```

```
4
```

```
N.relative_degree()
```

```
2
```

```
P.degree()
```

```
8
```

```
M.degree()*N.relative_degree() == P.degree()
```

```
True
```

Algebraic Numbers

Corollary 21.24 says that the set of *all* algebraic numbers forms a field. This field is implemented in Sage as `QQbar`. This allows for finding roots of polynomials as exact quantities which display as inexact numbers.

```
x = polygen(QQ, 'x')
p = x^4 + x^2 - 1
r = p.roots(ring=QQbar); r
```

```
[(-0.7861513777574233?,  1), (0.7861513777574233?,  1),
 (-1.272019649514069?*I, 1), (1.272019649514069?*I, 1)]
```

So we asked for the roots of a polynomial over the rationals, but requested any root that may lie outside the rationals and within the field of algebraic numbers. Since the field of algebraic numbers contains all such roots, we get a full four roots of the fourth-degree polynomial. These roots are computed to lie within an interval and the question mark

indicates that the preceding digits are accurate. (The integers paired with each root are the multiplicities of that root. Use the keyword `multiplicities=False` to turn them off.) Let us take a look under the hood and see how Sage manages the field of algebraic numbers.

```
r1 = r[0][0]; r1
```

-0.7861513777574233?

```
r1.as_number_field_element()
```

```
(Number Field in a with defining polynomial y^4 - y^2 - 1,
 a^3 - a,
 Ring morphism:
   From: Number Field in a with defining polynomial y^4 - y^2 - 1
   To:   Algebraic Real Field
   Defn: a |--> -1.272019649514069?)
```

Three items are associated with this initial root. First is a number field, with generator `a` and a defining polynomial similar to the polynomial we are finding the roots of, but not identical. Second is an expression in the generator `a`, which is the actual root. You might evaluate this expression with the numerical approximation of `a`, coming next, to verify this is a root. Finally, there is a ring homomorphism from the number field to the "Algebraic Real Field", `AA`, the subfield of `QQbar` with just real elements, which associates the generator `a` with the number `-1.272019649514069?`. Let us verify, in two ways, that the root given is really a root.

```
r1^4 + r1^2 - 1
```

0

```
N, rexact, homomorphism = r1.as_number_field_element()
(rexact)^4 + rexact^2 - 1
```

0

Now that we have enough theory to understand the field of algebraic numbers, and a natural way to represent them exactly, you might consider the operations in the field. If we take two algebraic numbers and add them together, we get another algebraic number (Corollary 21.24). So what is the resulting minimal polynomial? How is it computed in Sage? You could read the source code if you wanted the answer.

Geometric Constructions

Sage can do a lot of things, but it is not yet able to lay out lines with a straightedge and compass. However, we can very quickly determine that trisecting a 60 degree angle is impossible. We adjoin the cosine of a 20 degree angle (in radians) to the rationals, determine the degree of the extension, and check that it is not an integer power of 2. In one line. Sweet.

```
log(QQ[cos(pi/9)].degree(), 2) in ZZ
```

False

21.6   Sage Exercises

1.   Create the polynomial $p(x) = x^5 + 2x^4 + 1$ over $\mathbb{Z}_3$. Verify that it does not have any linear factors by evaluating $p(x)$ with each element of $\mathbb{Z}_3$, and then check that $p(x)$ is irreducible.

Create a finite field of order 3^5 with the `FiniteField()` command, but include the `modulus` keyword set to the polynomial $p(x)$ to override the default choice.

Recreate $p(x)$ as a polynomial over this field. Check each of the $3^5 = 243$ elements of the field to see if they are roots of the polynomial and list all of the elements which are roots. Finally, request that Sage give a factorization of $p(x)$ over the field, and comment on the relationship between your list of roots and your factorization.

2.   This problem continues the previous one. Build the ring of polynomials over $\mathbb{Z}_3$ and within this ring use $p(x)$ to generate a principal ideal. Finally construct the quotient of the polynomial ring by the ideal. Since the polynomial is irreducible, this quotient ring is a field, and by Proposition 21.12 this quotient ring is isomorphic to the number field in the previous problem.

Borrowing from your results in the previous question, construct five roots of the polynomial $p(x)$ within this quotient ring, but now as expressions in the generator of the quotient ring (which is technically a coset). Use Sage to verify that they are indeed roots. This demonstrates using a quotient ring to create a splitting field for an irreducible polynomial over a finite field.

3.   The subsection Algebraic Elements relies on techniques from linear algebra and contains Theorem 21.15: every finite extension is an algebraic extension. This exercise will help you understand this proof.

The polynomial $r(x) = x^4 + 2x + 2$ is irreducible over the rationals (Eisenstein's criterion with prime $p = 2$). Create a number field that contains a root of $r(x)$. By Theorem 21.15, and the remark following, every element of this finite field extension is an algebraic number, and hence satisfies some polynomial over the base field (it is this polynomial that Sage will produce with the `.minpoly()` method). This exercise will show how we can use just linear algebra to determine this minimal polynomial.

Suppose that `a` is the generator of the number field you just created with $r(x)$. Then we will determine the minimal polynomial of `t = 3a + 1` using just linear algebra. According to the proof, the first five powers of `t` (start counting from zero) will be linearly dependent. (Why?) So a nontrivial relation of linear dependence on these powers will provide the coefficients of a polynomial with `t` as a root. Compute these five powers, then construct the correct linear system to determine the coefficients of the minimal polynomial, solve the system, and suitably interpret its solutions.

Hints: The `vector()` and `matrix()` commands will create vectors and matrices, and the `.solve_right()` method for matrices can be used to find solutions. Given an element of the number field, which will necessarily be a polynomial in the generator `a`, the `.vector()` method of the element will provide the coefficients of this polynomial in a list.

4.   Construct the splitting field of $s(x) = x^4 + x^2 + 1$ and find a factorization of $s(x)$ over this field into linear factors.

5.   Form the number field, K, which contains a root of the irreducible polynomial $q(x) = x^3 + 3x^2 + 3x - 2$. Name your root `a`. Verify that $q(x)$ factors, but does not split, over K. With K now as the base field, form an extension of K where the quadratic factor of $q(x)$ has a root. Name this root `b`, and call this second extension of the tower L.

Use `M.<c> = L.absolute_field()` to form the flattened tower that is the absolute number field M. Find the defining polynomial of M with the `.polynomial()` method. From this

polynomial, which must have the generator `c` as a root, you should be able to use elementary algebra to write the generator as a fairly simple expression.

M should be the splitting field of $q(x)$. To see this, start over, and build from scratch a new number field, P, using the simple expression for `c` that you just found. Use `d` as the name of the root used to construct `P`. Since `d` is a root of the simple minimal polynomial for `c`, you should be able to write an expression for `d` that a pre-calculus student would recognize.

Now factor the original polynomial $q(x)$ (with rational coefficients) over P, to see the polynomial split (as expected). Using this factorization, and your simple expression for `d` write simplified expressions for the three roots of $q(x)$. See if you can convert between the two versions of the roots "by hand", and without using the isomorphisms provided by the `.structure()` method on `M`.

$$\boxed{22}$$

Finite Fields

Finite fields appear in many applications of algebra, including coding theory and cryptography. We already know one finite field, $\mathbb{Z}_p$, where p is prime. In this chapter we will show that a unique finite field of order p^n exists for every prime p, where n is a positive integer. Finite fields are also called Galois fields in honor of Évariste Galois, who was one of the first mathematicians to investigate them.

22.1 Structure of a Finite Field

Recall that a field F has **characteristic** p if p is the smallest positive integer such that for every nonzero element α in F, we have $p\alpha = 0$. If no such integer exists, then F has characteristic 0. From Theorem 16.19 we know that p must be prime. Suppose that F is a finite field with n elements. Then $n\alpha = 0$ for all α in F. Consequently, the characteristic of F must be p, where p is a prime dividing n. This discussion is summarized in the following proposition.

Proposition 22.1. *If F is a finite field, then the characteristic of F is p, where p is prime.*

Throughout this chapter we will assume that p is a prime number unless otherwise stated.

Proposition 22.2. *If F is a finite field of characteristic p, then the order of F is p^n for some $n \in \mathbb{N}$.*

PROOF. Let $\phi : \mathbb{Z} \to F$ be the ring homomorphism defined by $\phi(n) = n \cdot 1$. Since the characteristic of F is p, the kernel of ϕ must be $p\mathbb{Z}$ and the image of ϕ must be a subfield of F isomorphic to $\mathbb{Z}_p$. We will denote this subfield by K. Since F is a finite field, it must be a finite extension of K and, therefore, an algebraic extension of K. Suppose that $[F : K] = n$ is the dimension of F, where F is a K vector space. There must exist elements $\alpha_1, \ldots, \alpha_n \in F$ such that any element α in F can be written uniquely in the form

$$\alpha = a_1\alpha_1 + \cdots + a_n\alpha_n,$$

where the a_i's are in K. Since there are p elements in K, there are p^n possible linear combinations of the α_i's. Therefore, the order of F must be p^n. $\square$

Lemma 22.3 Freshman's Dream. *Let p be prime and D be an integral domain of characteristic p. Then*

$$a^{p^n} + b^{p^n} = (a + b)^{p^n}$$

for all positive integers n.

PROOF. We will prove this lemma using mathematical induction on n. We can use the binomial formula (see Chapter 2, Example 2.4) to verify the case for $n = 1$; that is,

$$(a + b)^p = \sum_{k=0}^{p} \binom{p}{k} a^k b^{p-k}.$$

If $0 < k < p$, then

$$\binom{p}{k} = \frac{p!}{k!(p-k)!}$$

must be divisible by p, since p cannot divide $k!(p-k)!$. Note that D is an integral domain of characteristic p, so all but the first and last terms in the sum must be zero. Therefore, $(a + b)^p = a^p + b^p$.

Now suppose that the result holds for all k, where $1 \leq k \leq n$. By the induction hypothesis,

$$(a + b)^{p^{n+1}} = ((a + b)^p)^{p^n} = (a^p + b^p)^{p^n} = (a^p)^{p^n} + (b^p)^{p^n} = a^{p^{n+1}} + b^{p^{n+1}}.$$

Therefore, the lemma is true for $n + 1$ and the proof is complete. $\square$

Let F be a field. A polynomial $f(x) \in F[x]$ of degree n is **separable** if it has n distinct roots in the splitting field of $f(x)$; that is, $f(x)$ is separable when it factors into distinct linear factors over the splitting field of f. An extension E of F is a **separable extension** of F if every element in E is the root of a separable polynomial in $F[x]$.

Example 22.4. The polynomial $x^2 - 2$ is separable over $\mathbb{Q}$ since it factors as $(x - \sqrt{2})(x + \sqrt{2})$. In fact, $\mathbb{Q}(\sqrt{2})$ is a separable extension of $\mathbb{Q}$. Let $\alpha = a + b\sqrt{2}$ be any element in $\mathbb{Q}(\sqrt{2})$. If $b = 0$, then α is a root of $x - a$. If $b \neq 0$, then α is the root of the separable polynomial

$$x^2 - 2ax + a^2 - 2b^2 = (x - (a + b\sqrt{2}))(x - (a - b\sqrt{2})).$$

Fortunately, we have an easy test to determine the separability of any polynomial. Let

$$f(x) = a_0 + a_1 x + \cdots + a_n x^n$$

be any polynomial in $F[x]$. Define the **derivative** of $f(x)$ to be

$$f'(x) = a_1 + 2a_2 x + \cdots + na_n x^{n-1}.$$

Lemma 22.5. *Let F be a field and $f(x) \in F[x]$. Then $f(x)$ is separable if and only if $f(x)$ and $f'(x)$ are relatively prime.*

PROOF. Let $f(x)$ be separable. Then $f(x)$ factors over some extension field of F as $f(x) = (x - \alpha_1)(x - \alpha_2) \cdots (x - \alpha_n)$, where $\alpha_i \neq \alpha_j$ for $i \neq j$. Taking the derivative of $f(x)$, we see that

$$\begin{aligned}
f'(x) &= (x - \alpha_2) \cdots (x - \alpha_n) \\
&+ (x - \alpha_1)(x - \alpha_3) \cdots (x - \alpha_n) \\
&+ \cdots + (x - \alpha_1) \cdots (x - \alpha_{n-1}).
\end{aligned}$$

Hence, $f(x)$ and $f'(x)$ can have no common factors.

To prove the converse, we will show that the contrapositive of the statement is true. Suppose that $f(x) = (x - \alpha)^k g(x)$, where $k > 1$. Differentiating, we have

$$f'(x) = k(x - \alpha)^{k-1} g(x) + (x - \alpha)^k g'(x).$$

Therefore, $f(x)$ and $f'(x)$ have a common factor. $\square$

Theorem 22.6. *For every prime p and every positive integer n, there exists a finite field F with p^n elements. Furthermore, any field of order p^n is isomorphic to the splitting field of $x^{p^n} - x$ over $\mathbb{Z}_p$.*

PROOF.   Let $f(x) = x^{p^n} - x$ and let F be the splitting field of $f(x)$. Then by Lemma 22.5, $f(x)$ has p^n distinct zeros in F, since $f'(x) = p^n x^{p^n - 1} - 1 = -1$ is relatively prime to $f(x)$. We claim that the roots of $f(x)$ form a subfield of F. Certainly 0 and 1 are zeros of $f(x)$. If α and β are zeros of $f(x)$, then $\alpha + \beta$ and $\alpha\beta$ are also zeros of $f(x)$, since $\alpha^{p^n} + \beta^{p^n} = (\alpha + \beta)^{p^n}$ and $\alpha^{p^n} \beta^{p^n} = (\alpha\beta)^{p^n}$. We also need to show that the additive inverse and the multiplicative inverse of each root of $f(x)$ are roots of $f(x)$. For any zero α of $f(x)$, we know that $-\alpha$ is also a zero of $f(x)$, since

$$f(-\alpha) = (-\alpha)^{p^n} - (-\alpha) = -\alpha^{p^n} + \alpha = -(\alpha^{p^n} - \alpha) = 0,$$

provided p is odd. If $p = 2$, then

$$f(-\alpha) = (-\alpha)^{2^n} - (-\alpha) = \alpha + \alpha = 0.$$

If $\alpha \neq 0$, then $(\alpha^{-1})^{p^n} = (\alpha^{p^n})^{-1} = \alpha^{-1}$. Since the zeros of $f(x)$ form a subfield of F and $f(x)$ splits in this subfield, the subfield must be all of F.

Let E be any other field of order p^n. To show that E is isomorphic to F, we must show that every element in E is a root of $f(x)$. Certainly 0 is a root of $f(x)$. Let α be a nonzero element of E. The order of the multiplicative group of nonzero elements of E is $p^n - 1$; hence, $\alpha^{p^n - 1} = 1$ or $\alpha^{p^n} - \alpha = 0$. Since E contains p^n elements, E must be a splitting field of $f(x)$; however, by Corollary 21.34, the splitting field of any polynomial is unique up to isomorphism. $\square$

The unique finite field with p^n elements is called the ***Galois field*** of order p^n. We will denote this field by $\mathrm{GF}(p^n)$.

Theorem 22.7. *Every subfield of the Galois field $\mathrm{GF}(p^n)$ has p^m elements, where m divides n. Conversely, if $m \mid n$ for $m > 0$, then there exists a unique subfield of $\mathrm{GF}(p^n)$ isomorphic to $\mathrm{GF}(p^m)$.*

PROOF.   Let F be a subfield of $E = \mathrm{GF}(p^n)$. Then F must be a field extension of K that contains p^m elements, where K is isomorphic to $\mathbb{Z}_p$. Then $m \mid n$, since $[E : K] = [E : F][F : K]$.

To prove the converse, suppose that $m \mid n$ for some $m > 0$. Then $p^m - 1$ divides $p^n - 1$. Consequently, $x^{p^m - 1} - 1$ divides $x^{p^n - 1} - 1$. Therefore, $x^{p^m} - x$ must divide $x^{p^n} - x$, and every zero of $x^{p^m} - x$ is also a zero of $x^{p^n} - x$. Thus, $\mathrm{GF}(p^n)$ contains, as a subfield, a splitting field of $x^{p^m} - x$, which must be isomorphic to $\mathrm{GF}(p^m)$. $\square$

Example 22.8. The lattice of subfields of $\mathrm{GF}(p^{24})$ is given in Figure 22.9.
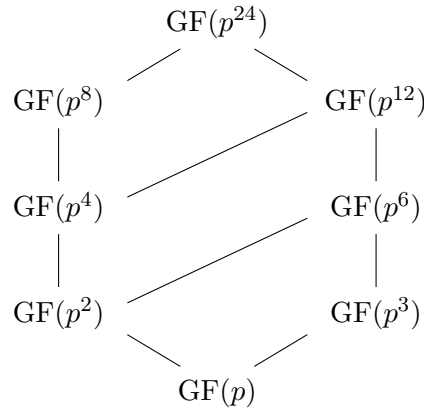
Figure 22.9: Subfields of $\mathrm{GF}(p^{24})$

With each field F we have a multiplicative group of nonzero elements of F which we will denote by F^*. The multiplicative group of any finite field is cyclic. This result follows from the more general result that we will prove in the next theorem.

Theorem 22.10. *If G is a finite subgroup of F^*, the multiplicative group of nonzero elements of a field F, then G is cyclic.*

PROOF. Let G be a finite subgroup of F^* of order n. By the Fundamental Theorem of Finite Abelian Groups (Theorem 13.4),

$$G \cong \mathbb{Z}_{p_1^{e_1}} \times \cdots \times \mathbb{Z}_{p_k^{e_k}},$$

where $n = p_1^{e_1} \cdots p_k^{e_k}$ and the $p_1, \ldots, p_k$ are (not necessarily distinct) primes. Let m be the least common multiple of $p_1^{e_1}, \ldots, p_k^{e_k}$. Then G contains an element of order m. Since every α in G satisfies $x^r - 1$ for some r dividing m, α must also be a root of $x^m - 1$. Since $x^m - 1$ has at most m roots in F, $n \leq m$. On the other hand, we know that $m \leq |G|$; therefore, $m = n$. Thus, G contains an element of order n and must be cyclic. $\square$

Corollary 22.11. *The multiplicative group of all nonzero elements of a finite field is cyclic.*

Corollary 22.12. *Every finite extension E of a finite field F is a simple extension of F.*

PROOF. Let α be a generator for the cyclic group E^* of nonzero elements of E. Then $E = F(\alpha)$. $\square$

Example 22.13. The finite field $\mathrm{GF}(2^4)$ is isomorphic to the field $\mathbb{Z}_2/\langle 1 + x + x^4 \rangle$. Therefore, the elements of $\mathrm{GF}(2^4)$ can be taken to be

$$\{a_0 + a_1\alpha + a_2\alpha^2 + a_3\alpha^3 : a_i \in \mathbb{Z}_2 \text{ and } 1 + \alpha + \alpha^4 = 0\}.$$

Remembering that $1 + \alpha + \alpha^4 = 0$, we add and multiply elements of $\mathrm{GF}(2^4)$ exactly as we add and multiply polynomials. The multiplicative group of $\mathrm{GF}(2^4)$ is isomorphic to $\mathbb{Z}_{15}$ with generator α:

$$
\begin{array}{lll}
\alpha^1 = \alpha & \alpha^6 = \alpha^2 + \alpha^3 & \alpha^{11} = \alpha + \alpha^2 + \alpha^3 \\
\alpha^2 = \alpha^2 & \alpha^7 = 1 + \alpha + \alpha^3 & \alpha^{12} = 1 + \alpha + \alpha^2 + \alpha^3 \\
\alpha^3 = \alpha^3 & \alpha^8 = 1 + \alpha^2 & \alpha^{13} = 1 + \alpha^2 + \alpha^3 \\
\alpha^4 = 1 + \alpha & \alpha^9 = \alpha + \alpha^3 & \alpha^{14} = 1 + \alpha^3 \\
\alpha^5 = \alpha + \alpha^2 & \alpha^{10} = 1 + \alpha + \alpha^2 & \alpha^{15} = 1.
\end{array}
$$

22.2   Polynomial Codes

With knowledge of polynomial rings and finite fields, it is now possible to derive more sophisticated codes than those of Chapter 8. First let us recall that an (n, k)-block code consists of a one-to-one encoding function $E : \mathbb{Z}_2^k \to \mathbb{Z}_2^n$ and a decoding function $D : \mathbb{Z}_2^n \to \mathbb{Z}_2^k$. The code is error-correcting if D is onto. A code is a linear code if it is the null space of a matrix $H \in \mathbb{M}_{k \times n}(\mathbb{Z}_2)$.

We are interested in a class of codes known as cyclic codes. Let $\phi : \mathbb{Z}_2^k \to \mathbb{Z}_2^n$ be a binary (n, k)-block code. Then ϕ is a ***cyclic code*** if for every codeword $(a_1, a_2, \ldots, a_n)$, the cyclically shifted n-tuple $(a_n, a_1, a_2, \ldots, a_{n-1})$ is also a codeword. Cyclic codes are particularly easy to implement on a computer using shift registers [2, 3].

Example 22.14. Consider the $(6, 3)$-linear codes generated by the two matrices

$$
G_1 = \begin{pmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1 \\ 1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1 \end{pmatrix} \quad \text{and} \quad G_2 = \begin{pmatrix} 1 & 0 & 0 \\ 1 & 1 & 0 \\ 1 & 1 & 1 \\ 1 & 1 & 1 \\ 0 & 1 & 1 \\ 0 & 0 & 1 \end{pmatrix}.
$$

Messages in the first code are encoded as follows:

$$
\begin{array}{llll}
(000) & \mapsto & (000000) & \qquad (100) \mapsto (100100) \\
(001) & \mapsto & (001001) & \qquad (101) \mapsto (101101) \\
(010) & \mapsto & (010010) & \qquad (110) \mapsto (110110) \\
(011) & \mapsto & (011011) & \qquad (111) \mapsto (111111).
\end{array}
$$

It is easy to see that the codewords form a cyclic code. In the second code, 3-tuples are encoded in the following manner:

$$
\begin{array}{llll}
(000) & \mapsto & (000000) & \qquad (100) \mapsto (111100) \\
(001) & \mapsto & (001111) & \qquad (101) \mapsto (110011) \\
(010) & \mapsto & (011110) & \qquad (110) \mapsto (100010) \\
(011) & \mapsto & (010001) & \qquad (111) \mapsto (101101).
\end{array}
$$

This code cannot be cyclic, since (101101) is a codeword but (011011) is not a codeword.

Polynomial Codes

We would like to find an easy method of obtaining cyclic linear codes. To accomplish this, we can use our knowledge of finite fields and polynomial rings over $\mathbb{Z}_2$. Any binary n-tuple can be interpreted as a polynomial in $\mathbb{Z}_2[x]$. Stated another way, the n-tuple $(a_0, a_1, \ldots, a_{n-1})$ corresponds to the polynomial

$$
f(x) = a_0 + a_1 x + \cdots + a_{n-1} x^{n-1},
$$

where the degree of $f(x)$ is at most $n - 1$. For example, the polynomial corresponding to the 5-tuple (10011) is

$$
1 + 0x + 0x^2 + 1x^3 + 1x^4 = 1 + x^3 + x^4.
$$

Conversely, with any polynomial $f(x) \in \mathbb{Z}_2[x]$ with $\deg f(x) < n$ we can associate a binary n-tuple. The polynomial $x + x^2 + x^4$ corresponds to the 5-tuple (01101).

Let us fix a nonconstant polynomial $g(x)$ in $\mathbb{Z}_2[x]$ of degree $n - k$. We can define an (n, k)-code C in the following manner. If $(a_0, \ldots, a_{k-1})$ is a k-tuple to be encoded, then $f(x) = a_0 + a_1 x + \cdots + a_{k-1} x^{k-1}$ is the corresponding polynomial in $\mathbb{Z}_2[x]$. To encode $f(x)$, we multiply by $g(x)$. The codewords in C are all those polynomials in $\mathbb{Z}_2[x]$ of degree less than n that are divisible by $g(x)$. Codes obtained in this manner are called **polynomial codes**.

Example 22.15. If we let $g(x) = 1 + x^3$, we can define a $(6, 3)$-code C as follows. To encode a 3-tuple (a_0, a_1, a_2), we multiply the corresponding polynomial $f(x) = a_0 + a_1 x + a_2 x^2$ by $1 + x^3$. We are defining a map $\phi : \mathbb{Z}_2^3 \to \mathbb{Z}_2^6$ by $\phi : f(x) \mapsto g(x) f(x)$. It is easy to check that this map is a group homomorphism. In fact, if we regard $\mathbb{Z}_2^n$ as a vector space over $\mathbb{Z}_2$, ϕ is a linear transformation of vector spaces (see Exercise 20.4.15, Chapter 20). Let us compute the kernel of ϕ. Observe that $\phi(a_0, a_1, a_2) = (000000)$ exactly when

$$0 + 0x + 0x^2 + 0x^3 + 0x^4 + 0x^5 = (1 + x^3)(a_0 + a_1 x + a_2 x^2)$$
$$= a_0 + a_1 x + a_2 x^2 + a_0 x^3 + a_1 x^4 + a_2 x^5.$$

Since the polynomials over a field form an integral domain, $a_0 + a_1 x + a_2 x^2$ must be the zero polynomial. Therefore, $\ker \phi = \{(000)\}$ and ϕ is one-to-one.

To calculate a generator matrix for C, we merely need to examine the way the polynomials 1, x, and x^2 are encoded:

$$(1 + x^3) \cdot 1 = 1 + x^3$$
$$(1 + x^3)x = x + x^4$$
$$(1 + x^3)x^2 = x^2 + x^5.$$

We obtain the code corresponding to the generator matrix G_1 in Example 22.14. The parity-check matrix for this code is

$$H = \begin{pmatrix} 1 & 0 & 0 & 1 & 0 & 0 \\ 0 & 1 & 0 & 0 & 1 & 0 \\ 0 & 0 & 1 & 0 & 0 & 1 \end{pmatrix}.$$

Since the smallest weight of any nonzero codeword is 2, this code has the ability to detect all single errors.

Rings of polynomials have a great deal of structure; therefore, our immediate goal is to establish a link between polynomial codes and ring theory. Recall that $x^n - 1 = (x - 1)(x^{n-1} + \cdots + x + 1)$. The factor ring

$$R_n = \mathbb{Z}_2[x]/\langle x^n - 1 \rangle$$

can be considered to be the ring of polynomials of the form

$$f(t) = a_0 + a_1 t + \cdots + a_{n-1} t^{n-1}$$

that satisfy the condition $t^n = 1$. It is an easy exercise to show that $\mathbb{Z}_2^n$ and R_n are isomorphic as vector spaces. We will often identify elements in $\mathbb{Z}_2^n$ with elements in $\mathbb{Z}[x]/\langle x^n - 1 \rangle$. In this manner we can interpret a linear code as a subset of $\mathbb{Z}[x]/\langle x^n - 1 \rangle$.

The additional ring structure on polynomial codes is very powerful in describing cyclic codes. A cyclic shift of an n-tuple can be described by polynomial multiplication. If $f(t) = a_0 + a_1 t + \cdots + a_{n-1} t^{n-1}$ is a code polynomial in R_n, then

$$t f(t) = a_{n-1} + a_0 t + \cdots + a_{n-2} t^{n-1}$$

is the cyclically shifted word obtained from multiplying $f(t)$ by t. The following theorem gives a beautiful classification of cyclic codes in terms of the ideals of R_n.

Theorem 22.16. *A linear code C in $\mathbb{Z}_2^n$ is cyclic if and only if it is an ideal in $R_n = \mathbb{Z}[x]/\langle x^n - 1 \rangle$.*

PROOF.   Let C be a linear cyclic code and suppose that $f(t)$ is in C. Then $tf(t)$ must also be in C. Consequently, $t^k f(t)$ is in C for all $k \in \mathbb{N}$. Since C is a linear code, any linear combination of the codewords $f(t), tf(t), t^2 f(t), \dots, t^{n-1} f(t)$ is also a codeword; therefore, for every polynomial $p(t)$, $p(t)f(t)$ is in C. Hence, C is an ideal.

Conversely, let C be an ideal in $\mathbb{Z}_2[x]/\langle x^n + 1 \rangle$. Suppose that $f(t) = a_0 + a_1 t + \cdots + a_{n-1} t^{n-1}$ is a codeword in C. Then $tf(t)$ is a codeword in C; that is, $(a_1, \dots, a_{n-1}, a_0)$ is in C. $\square$

Theorem 22.16 tells us that knowing the ideals of R_n is equivalent to knowing the linear cyclic codes in $\mathbb{Z}_2^n$. Fortunately, the ideals in R_n are easy to describe. The natural ring homomorphism $\phi : \mathbb{Z}_2[x] \to R_n$ defined by $\phi[f(x)] = f(t)$ is a surjective homomorphism. The kernel of ϕ is the ideal generated by $x^n - 1$. By Theorem 16.34, every ideal C in R_n is of the form $\phi(I)$, where I is an ideal in $\mathbb{Z}_2[x]$ that contains $\langle x^n - 1 \rangle$. By Theorem 17.20, we know that every ideal I in $\mathbb{Z}_2[x]$ is a principal ideal, since $\mathbb{Z}_2$ is a field. Therefore, $I = \langle g(x) \rangle$ for some unique monic polynomial in $\mathbb{Z}_2[x]$. Since $\langle x^n - 1 \rangle$ is contained in I, it must be the case that $g(x)$ divides $x^n - 1$. Consequently, every ideal C in R_n is of the form

$$C = \langle g(t) \rangle = \{ f(t)g(t) : f(t) \in R_n \text{ and } g(x) \mid (x^n - 1) \text{ in } \mathbb{Z}_2[x] \}.$$

The unique monic polynomial of the smallest degree that generates C is called the **minimal generator polynomial** of C.

Example 22.17. If we factor $x^7 - 1$ into irreducible components, we have

$$x^7 - 1 = (1 + x)(1 + x + x^3)(1 + x^2 + x^3).$$

We see that $g(t) = (1 + t + t^3)$ generates an ideal C in R_7. This code is a $(7, 4)$-block code. As in Example 22.15, it is easy to calculate a generator matrix by examining what $g(t)$ does to the polynomials 1, t, t^2, and t^3. A generator matrix for C is

$$G = \begin{pmatrix} 1 & 0 & 0 & 0 \\ 1 & 1 & 0 & 0 \\ 0 & 1 & 1 & 0 \\ 1 & 0 & 1 & 1 \\ 0 & 1 & 0 & 1 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1 \end{pmatrix}.$$

In general, we can determine a generator matrix for an (n, k)-code C by the manner in which the elements t^k are encoded. Let $x^n - 1 = g(x)h(x)$ in $\mathbb{Z}_2[x]$. If $g(x) = g_0 + g_1 x + \cdots + g_{n-k} x^{n-k}$ and $h(x) = h_0 + h_1 x + \cdots + h_k x^k$, then the $n \times k$ matrix

$$G = \begin{pmatrix} g_0 & 0 & \cdots & 0 \\ g_1 & g_0 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ g_{n-k} & g_{n-k-1} & \cdots & g_0 \\ 0 & g_{n-k} & \cdots & g_1 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & g_{n-k} \end{pmatrix}$$

is a generator matrix for the code C with generator polynomial $g(t)$. The parity-check matrix for C is the $(n - k) \times n$ matrix

$$H = \begin{pmatrix} 0 & \cdots & 0 & 0 & h_k & \cdots & h_0 \\ 0 & \cdots & 0 & h_k & \cdots & h_0 & 0 \\ \cdots & \cdots & \cdots & \cdots & \cdots & \cdots & \cdots \\ h_k & \cdots & h_0 & 0 & 0 & \cdots & 0 \end{pmatrix}.$$

We will leave the details of the proof of the following proposition as an exercise.

Proposition 22.18. *Let* $C = \langle g(t) \rangle$ *be a cyclic code in* R_n *and suppose that* $x^n - 1 = g(x)h(x)$. *Then* G *and* H *are generator and parity-check matrices for* C, *respectively. Furthermore,* $HG = 0$.

Example 22.19. In Example 22.17,

$$x^7 - 1 = g(x)h(x) = (1 + x + x^3)(1 + x + x^2 + x^4).$$

Therefore, a parity-check matrix for this code is

$$H = \begin{pmatrix} 0 & 0 & 1 & 0 & 1 & 1 & 1 \\ 0 & 1 & 0 & 1 & 1 & 1 & 0 \\ 1 & 0 & 1 & 1 & 1 & 0 & 0 \end{pmatrix}.$$

To determine the error-detecting and error-correcting capabilities of a cyclic code, we need to know something about determinants. If $\alpha_1, \ldots, \alpha_n$ are elements in a field F, then the $n \times n$ matrix

$$\begin{pmatrix} 1 & 1 & \cdots & 1 \\ \alpha_1 & \alpha_2 & \cdots & \alpha_n \\ \alpha_1^2 & \alpha_2^2 & \cdots & \alpha_n^2 \\ \vdots & \vdots & \ddots & \vdots \\ \alpha_1^{n-1} & \alpha_2^{n-1} & \cdots & \alpha_n^{n-1} \end{pmatrix}$$

is called the ***Vandermonde matrix***. The determinant of this matrix is called the ***Vandermonde determinant***. We will need the following lemma in our investigation of cyclic codes.

Lemma 22.20. *Let* $\alpha_1, \ldots, \alpha_n$ *be elements in a field* F *with* $n \geq 2$. *Then*

$$\det \begin{pmatrix} 1 & 1 & \cdots & 1 \\ \alpha_1 & \alpha_2 & \cdots & \alpha_n \\ \alpha_1^2 & \alpha_2^2 & \cdots & \alpha_n^2 \\ \vdots & \vdots & \ddots & \vdots \\ \alpha_1^{n-1} & \alpha_2^{n-1} & \cdots & \alpha_n^{n-1} \end{pmatrix} = \prod_{1 \leq j < i \leq n} (\alpha_i - \alpha_j).$$

In particular, if the α_i*'s are distinct, then the determinant is nonzero.*

PROOF. We will induct on n. If $n = 2$, then the determinant is $\alpha_2 - \alpha_1$. Let us assume the result for $n - 1$ and consider the polynomial $p(x)$ defined by

$$p(x) = \det \begin{pmatrix} 1 & 1 & \cdots & 1 & 1 \\ \alpha_1 & \alpha_2 & \cdots & \alpha_{n-1} & x \\ \alpha_1^2 & \alpha_2^2 & \cdots & \alpha_{n-1}^2 & x^2 \\ \vdots & \vdots & \ddots & \vdots & \vdots \\ \alpha_1^{n-1} & \alpha_2^{n-1} & \cdots & \alpha_{n-1}^{n-1} & x^{n-1} \end{pmatrix}.$$

Expanding this determinant by cofactors on the last column, we see that $p(x)$ is a polynomial of at most degree $n-1$. Moreover, the roots of $p(x)$ are $\alpha_1, \ldots, \alpha_{n-1}$, since the substitution of any one of these elements in the last column will produce a column identical to the last column in the matrix. Remember that the determinant of a matrix is zero if it has two identical columns. Therefore,

$$p(x) = (x - \alpha_1)(x - \alpha_2) \cdots (x - \alpha_{n-1})\beta,$$

where

$$\beta = (-1)^{n+n} \det \begin{pmatrix} 1 & 1 & \cdots & 1 \\ \alpha_1 & \alpha_2 & \cdots & \alpha_{n-1} \\ \alpha_1^2 & \alpha_2^2 & \cdots & \alpha_{n-1}^2 \\ \vdots & \vdots & \ddots & \vdots \\ \alpha_1^{n-2} & \alpha_2^{n-2} & \cdots & \alpha_{n-1}^{n-2} \end{pmatrix}.$$

By our induction hypothesis,

$$\beta = (-1)^{n+n} \prod_{1 \le j < i \le n-1} (\alpha_i - \alpha_j).$$

If we let $x = \alpha_n$, the result now follows immediately. $\qquad\square$

The following theorem gives us an estimate on the error detection and correction capabilities for a particular generator polynomial.

Theorem 22.21. *Let $C = \langle g(t) \rangle$ be a cyclic code in R_n and suppose that ω is a primitive nth root of unity over $\mathbb{Z}_2$. If s consecutive powers of ω are roots of $g(x)$, then the minimum distance of C is at least $s + 1$.*

PROOF.   Suppose that

$$g(\omega^r) = g(\omega^{r+1}) = \cdots = g(\omega^{r+s-1}) = 0.$$

Let $f(x)$ be some polynomial in C with s or fewer nonzero coefficients. We can assume that

$$f(x) = a_{i_0} x^{i_0} + a_{i_1} x^{i_1} + \cdots + a_{i_{s-1}} x^{i_{s-1}}$$

be some polynomial in C. It will suffice to show that all of the a_i's must be 0. Since

$$g(\omega^r) = g(\omega^{r+1}) = \cdots = g(\omega^{r+s-1}) = 0$$

and $g(x)$ divides $f(x)$,

$$f(\omega^r) = f(\omega^{r+1}) = \cdots = f(\omega^{r+s-1}) = 0.$$

Equivalently, we have the following system of equations:

$$a_{i_0}(\omega^r)^{i_0} + a_{i_1}(\omega^r)^{i_1} + \cdots + a_{i_{s-1}}(\omega^r)^{i_{s-1}} = 0$$
$$a_{i_0}(\omega^{r+1})^{i_0} + a_{i_1}(\omega^{r+1})^{i_2} + \cdots + a_{i_{s-1}}(\omega^{r+1})^{i_{s-1}} = 0$$
$$\vdots$$
$$a_{i_0}(\omega^{r+s-1})^{i_0} + a_{i_1}(\omega^{r+s-1})^{i_1} + \cdots + a_{i_{s-1}}(\omega^{r+s-1})^{i_{s-1}} = 0.$$

Therefore, $(a_{i_0}, a_{i_1}, \ldots, a_{i_{s-1}})$ is a solution to the homogeneous system of linear equations

$$(\omega^{i_0})^r x_0 + (\omega^{i_1})^r x_1 + \cdots + (\omega^{i_{s-1}})^r x_{n-1} = 0$$

$$(\omega^{i_0})^{r+1}x_0 + (\omega^{i_1})^{r+1}x_1 + \cdots + (\omega^{i_{s-1}})^{r+1}x_{n-1} = 0$$

$$\vdots$$

$$(\omega^{i_0})^{r+s-1}x_0 + (\omega^{i_1})^{r+s-1}x_1 + \cdots + (\omega^{i_{s-1}})^{r+s-1}x_{n-1} = 0.$$

However, this system has a unique solution, since the determinant of the matrix

$$\begin{pmatrix} (\omega^{i_0})^r & (\omega^{i_1})^r & \cdots & (\omega^{i_{s-1}})^r \\ (\omega^{i_0})^{r+1} & (\omega^{i_1})^{r+1} & \cdots & (\omega^{i_{s-1}})^{r+1} \\ \vdots & \vdots & \ddots & \vdots \\ (\omega^{i_0})^{r+s-1} & (\omega^{i_1})^{r+s-1} & \cdots & (\omega^{i_{s-1}})^{r+s-1} \end{pmatrix}$$

can be shown to be nonzero using Lemma 22.20 and the basic properties of determinants (Exercise). Therefore, this solution must be $a_{i_0} = a_{i_1} = \cdots = a_{i_{s-1}} = 0$. $\square$

BCH Codes

Some of the most important codes, discovered independently by A. Hocquenghem in 1959 and by R. C. Bose and D. V. Ray-Chaudhuri in 1960, are BCH codes. The European and transatlantic communication systems both use BCH codes. Information words to be encoded are of length 231, and a polynomial of degree 24 is used to generate the code. Since $231 + 24 = 255 = 2^8 - 1$, we are dealing with a $(255, 231)$-block code. This BCH code will detect six errors and has a failure rate of 1 in 16 million. One advantage of BCH codes is that efficient error correction algorithms exist for them.

The idea behind BCH codes is to choose a generator polynomial of smallest degree that has the largest error detection and error correction capabilities. Let $d = 2r + 1$ for some $r \geq 0$. Suppose that ω is a primitive nth root of unity over $\mathbb{Z}_2$, and let $m_i(x)$ be the minimal polynomial over $\mathbb{Z}_2$ of ω^i. If

$$g(x) = \mathrm{lcm}[m_1(x), m_2(x), \ldots, m_{2r}(x)],$$

then the cyclic code $\langle g(t) \rangle$ in R_n is called the **bch** *code of length* n *and distance* d. By Theorem 22.21, the minimum distance of C is at least d.

Theorem 22.22. *Let* $C = \langle g(t) \rangle$ *be a cyclic code in* R_n. *The following statements are equivalent.*

1. *The code* C *is a* BCH *code whose minimum distance is at least* d.

2. *A code polynomial* $f(t)$ *is in* C *if and only if* $f(\omega^i) = 0$ *for* $1 \leq i < d$.

3. *The matrix*

$$H = \begin{pmatrix} 1 & \omega & \omega^2 & \cdots & \omega^{n-1} \\ 1 & \omega^2 & \omega^4 & \cdots & \omega^{(n-1)(2)} \\ 1 & \omega^3 & \omega^6 & \cdots & \omega^{(n-1)(3)} \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ 1 & \omega^{2r} & \omega^{4r} & \cdots & \omega^{(n-1)(2r)} \end{pmatrix}$$

is a parity-check matrix for C.

PROOF. $(1) \Rightarrow (2)$. If $f(t)$ is in C, then $g(x) \mid f(x)$ in $\mathbb{Z}_2[x]$. Hence, for $i = 1, \ldots, 2r$, $f(\omega^i) = 0$ since $g(\omega^i) = 0$. Conversely, suppose that $f(\omega^i) = 0$ for $1 \leq i \leq d$. Then $f(x)$ is

divisible by each $m_i(x)$, since $m_i(x)$ is the minimal polynomial of ω^i. Therefore, $g(x) \mid f(x)$ by the definition of $g(x)$. Consequently, $f(x)$ is a codeword.

$(2) \Rightarrow (3)$. Let $f(t) = a_0 + a_1 t + \cdots + a_{n-1} v t^{n-1}$ be in R_n. The corresponding n-tuple in $\mathbb{Z}_2^n$ is $\mathbf{x} = (a_0 a_1 \cdots a_{n-1})^t$. By (2),

$$
H\mathbf{x} = \begin{pmatrix} a_0 + a_1 \omega + \cdots + a_{n-1} \omega^{n-1} \\ a_0 + a_1 \omega^2 + \cdots + a_{n-1} (\omega^2)^{n-1} \\ \vdots \\ a_0 + a_1 \omega^{2r} + \cdots + a_{n-1} (\omega^{2r})^{n-1} \end{pmatrix} = \begin{pmatrix} f(\omega) \\ f(\omega^2) \\ \vdots \\ f(\omega^{2r}) \end{pmatrix} = 0
$$

exactly when $f(t)$ is in C. Thus, H is a parity-check matrix for C.

$(3) \Rightarrow (1)$. By (3), a code polynomial $f(t) = a_0 + a_1 t + \cdots + a_{n-1} t^{n-1}$ is in C exactly when $f(\omega^i) = 0$ for $i = 1, \ldots, 2r$. The smallest such polynomial is $g(t) = \mathrm{lcm}[m_1(t), \ldots, m_{2r}(t)]$. Therefore, $C = \langle g(t) \rangle$. $\square$

Example 22.23. It is easy to verify that $x^{15} - 1 \in \mathbb{Z}_2[x]$ has a factorization

$$
x^{15} - 1 = (x + 1)(x^2 + x + 1)(x^4 + x + 1)(x^4 + x^3 + 1)(x^4 + x^3 + x^2 + x + 1),
$$

where each of the factors is an irreducible polynomial. Let ω be a root of $1 + x + x^4$. The Galois field $\mathrm{GF}(2^4)$ is

$$
\{a_0 + a_1 \omega + a_2 \omega^2 + a_3 \omega^3 : a_i \in \mathbb{Z}_2 \text{ and } 1 + \omega + \omega^4 = 0\}.
$$

By Example 22.8, ω is a primitive 15th root of unity. The minimal polynomial of ω is $m_1(x) = 1 + x + x^4$. It is easy to see that ω^2 and ω^4 are also roots of $m_1(x)$. The minimal polynomial of ω^3 is $m_2(x) = 1 + x + x^2 + x^3 + x^4$. Therefore,

$$
g(x) = m_1(x) m_2(x) = 1 + x^4 + x^6 + x^7 + x^8
$$

has roots ω, ω^2, ω^3, ω^4. Since both $m_1(x)$ and $m_2(x)$ divide $x^{15} - 1$, the BCH code is a $(15, 7)$-code. If $x^{15} - 1 = g(x)h(x)$, then $h(x) = 1 + x^4 + x^6 + x^7$; therefore, a parity-check matrix for this code is

$$
\begin{pmatrix}
0 & 0 & 0 & 0 & 0 & 0 & 0 & 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 \\
0 & 0 & 0 & 0 & 0 & 0 & 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 & 0 \\
0 & 0 & 0 & 0 & 0 & 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 \\
0 & 0 & 0 & 0 & 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 & 0 \\
0 & 0 & 0 & 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 \\
0 & 0 & 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 \\
0 & 1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 \\
1 & 1 & 0 & 1 & 0 & 0 & 0 & 1 & 0 & 0 & 0 & 0 & 0 & 0 & 0
\end{pmatrix}.
$$

22.3   Exercises

1.   Calculate each of the following.

(a) $[\mathrm{GF}(3^6) : \mathrm{GF}(3^3)]$

(b) $[\mathrm{GF}(128) : \mathrm{GF}(16)]$

(c) $[\mathrm{GF}(625) : \mathrm{GF}(25)]$

(d) $[\mathrm{GF}(p^{12}) : \mathrm{GF}(p^2)]$

2.   Calculate $[\mathrm{GF}(p^m) : \mathrm{GF}(p^n)]$, where $n \mid m$.

3. What is the lattice of subfields for $\mathrm{GF}(p^{30})$?

4. Let α be a zero of $x^3 + x^2 + 1$ over $\mathbb{Z}_2$. Construct a finite field of order 8. Show that $x^3 + x^2 + 1$ splits in $\mathbb{Z}_2(\alpha)$.

5. Construct a finite field of order 27.

6. Prove or disprove: $\mathbb{Q}^*$ is cyclic.

7. Factor each of the following polynomials in $\mathbb{Z}_2[x]$.

(a) $x^5 - 1$ (c) $x^9 - 1$

(b) $x^6 + x^5 + x^4 + x^3 + x^2 + x + 1$ (d) $x^4 + x^3 + x^2 + x + 1$

8. Prove or disprove: $\mathbb{Z}_2[x]/\langle x^3 + x + 1\rangle \cong \mathbb{Z}_2[x]/\langle x^3 + x^2 + 1\rangle$.

9. Determine the number of cyclic codes of length n for $n = 6, 7, 8, 10$.

10. Prove that the ideal $\langle t + 1\rangle$ in R_n is the code in $\mathbb{Z}_2^n$ consisting of all words of even parity.

11. Construct all BCH codes of

(a) length 7.   (b) length 15.

12. Prove or disprove: There exists a finite field that is algebraically closed.

13. Let p be prime. Prove that the field of rational functions $\mathbb{Z}_p(x)$ is an infinite field of characteristic p.

14. Let D be an integral domain of characteristic p. Prove that $(a - b)^{p^n} = a^{p^n} - b^{p^n}$ for all $a, b \in D$.

15. Show that every element in a finite field can be written as the sum of two squares.

16. Let E and F be subfields of a finite field K. If E is isomorphic to F, show that $E = F$.

17. Let $F \subset E \subset K$ be fields. If K is a separable extension of F, show that K is also separable extension of E.

18. Let E be an extension of a finite field F, where F has q elements. Let $\alpha \in E$ be algebraic over F of degree n. Prove that $F(\alpha)$ has q^n elements.

19. Show that every finite extension of a finite field F is simple; that is, if E is a finite extension of a finite field F, prove that there exists an $\alpha \in E$ such that $E = F(\alpha)$.

20. Show that for every n there exists an irreducible polynomial of degree n in $\mathbb{Z}_p[x]$.

21. Prove that the **Frobenius map** $\Phi : \mathrm{GF}(p^n) \to \mathrm{GF}(p^n)$ given by $\Phi : \alpha \mapsto \alpha^p$ is an automorphism of order n.

22. Show that every element in $\mathrm{GF}(p^n)$ can be written in the form a^p for some unique $a \in \mathrm{GF}(p^n)$.

23. Let E and F be subfields of $\mathrm{GF}(p^n)$. If $|E| = p^r$ and $|F| = p^s$, what is the order of $E \cap F$?

24. Wilson's Theorem. Let p be prime. Prove that $(p - 1)! \equiv -1 \pmod{p}$.

25. If $g(t)$ is the minimal generator polynomial for a cyclic code C in R_n, prove that the constant term of $g(x)$ is 1.

26. Often it is conceivable that a burst of errors might occur during transmission, as in the case of a power surge. Such a momentary burst of interference might alter several consecutive bits in a codeword. Cyclic codes permit the detection of such error bursts. Let C be an (n, k)-cyclic code. Prove that any error burst up to $n - k$ digits can be detected.

27.   Prove that the rings R_n and $\mathbb{Z}_2^n$ are isomorphic as vector spaces.

28.   Let C be a code in R_n that is generated by $g(t)$. If $\langle f(t) \rangle$ is another code in R_n, show that $\langle g(t) \rangle \subset \langle f(t) \rangle$ if and only if $f(x)$ divides $g(x)$ in $\mathbb{Z}_2[x]$.

29.   Let $C = \langle g(t) \rangle$ be a cyclic code in R_n and suppose that $x^n - 1 = g(x)h(x)$, where $g(x) = g_0 + g_1 x + \cdots + g_{n-k} x^{n-k}$ and $h(x) = h_0 + h_1 x + \cdots + h_k x^k$. Define G to be the $n \times k$ matrix

$$
G = \begin{pmatrix}
g_0 & 0 & \cdots & 0 \\
g_1 & g_0 & \cdots & 0 \\
\vdots & \vdots & \ddots & \vdots \\
g_{n-k} & g_{n-k-1} & \cdots & g_0 \\
0 & g_{n-k} & \cdots & g_1 \\
\vdots & \vdots & \ddots & \vdots \\
0 & 0 & \cdots & g_{n-k}
\end{pmatrix}
$$

and H to be the $(n-k) \times n$ matrix

$$
H = \begin{pmatrix}
0 & \cdots & 0 & 0 & h_k & \cdots & h_0 \\
0 & \cdots & 0 & h_k & \cdots & h_0 & 0 \\
\cdots & \cdots & \cdots & \cdots & \cdots & \cdots & \cdots \\
h_k & \cdots & h_0 & 0 & 0 & \cdots & 0
\end{pmatrix}.
$$

(a) Prove that G is a generator matrix for C.

(b) Prove that H is a parity-check matrix for C.

(c) Show that $HG = 0$.

22.4   Additional Exercises: Error Correction for BCH Codes

BCH codes have very attractive error correction algorithms. Let C be a BCH code in R_n, and suppose that a code polynomial $c(t) = c_0 + c_1 t + \cdots + c_{n-1} t^{n-1}$ is transmitted. Let $w(t) = w_0 + w_1 t + \cdots w_{n-1} t^{n-1}$ be the polynomial in R_n that is received. If errors have occurred in bits $a_1, \ldots, a_k$, then $w(t) = c(t) + e(t)$, where $e(t) = t^{a_1} + t^{a_2} + \cdots + t^{a_k}$ is the **error polynomial**. The decoder must determine the integers a_i and then recover $c(t)$ from $w(t)$ by flipping the a_ith bit. From $w(t)$ we can compute $w(\omega^i) = s_i$ for $i = 1, \ldots, 2r$, where ω is a primitive nth root of unity over $\mathbb{Z}_2$. We say the **syndrome** of $w(t)$ is $s_1, \ldots, s_{2r}$.**1.** Show that $w(t)$ is a code polynomial if and only if $s_i = 0$ for all i.

2.   Show that

$$
s_i = w(\omega^i) = e(\omega^i) = \omega^{ia_1} + \omega^{ia_2} + \cdots + \omega^{ia_k}
$$

for $i = 1, \ldots, 2r$. The **error-locator polynomial** is defined to be

$$
s(x) = (x + \omega^{a_1})(x + \omega^{a_2}) \cdots (x + \omega^{a_k}).
$$

3.   Recall the $(15, 7)$-block BCH code in Example 22.19. By Theorem 8.13, this code is capable of correcting two errors. Suppose that these errors occur in bits a_1 and a_2. The error-locator polynomial is $s(x) = (x + \omega^{a_1})(x + \omega^{a_2})$. Show that

$$
s(x) = x^2 + s_1 x + \left(s_1^2 + \frac{s_3}{s_1} \right).
$$

4. Let $w(t) = 1 + t^2 + t^4 + t^5 + t^7 + t^{12} + t^{13}$. Determine what the originally transmitted code polynomial was.

References and Suggested Readings

[1] Childs, L. *A Concrete Introduction to Higher Algebra.* 2nd ed. Springer-Verlag, New York, 1995.

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[3] Lidl, R. and Pilz, G. *Applied Abstract Algebra.* 2nd ed. Springer, New York, 1998. An excellent presentation of finite fields and their applications.

[4] Mackiw, G. *Applications of Abstract Algebra.* Wiley, New York, 1985.

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[6] van Lint, J. H. *Introduction to Coding Theory.* Springer, New York, 1999.

22.5 Sage

You have noticed in this chapter that finite fields have a great deal of structure. We have also seen finite fields in Sage regularly as examples of rings and fields. Now we can combine the two, mostly using commands we already know, plus a few new ones.

Creating Finite Fields

By Theorem 22.6 we know that all finite fields of a given order are isomorphic and that possible orders are limited to powers of primes. We can use the `FiniteField()` command, as before, or a shorter equivalent is `GF()`. Optionally, we can specify an irreducible polynomial for the contruction of the field. We can view this polynomial as the generator of the principal ideal of a polynomial ring, or we can view it as a "re-writing" rule for powers of the field's generator that allow us to multiply elements and reformulate them as linear combinations of lesser powers.

Absent providing an irreducible polynomial, Sage will use a Conway polynomial. You can determine these with the `conway_polynomial()` command, or just build a finite field and request the defining polynomial with the `.polynomial()` method.

```
F.<a> = GF(7^15); F
```

```
Finite Field in a of size 7^15
```

```
F.polynomial()
```

```
a^15 + 5*a^6 + 6*a^5 + 6*a^4 + 4*a^3 + a^2 + 2*a + 4
```

```
a^15 + 5*a^6 + 6*a^5 + 6*a^4 + 4*a^3 + a^2 + 2*a + 4
```

```
0
```

```
conway_polynomial(7, 15)
```

```
x^15 + 5*x^6 + 6*x^5 + 6*x^4 + 4*x^3 + x^2 + 2*x + 4
```

Just to be more readable, we coerce a list of coefficients into the set of polynomials (obtained with the `.parent()` method on a simple polynomial) to define a polynomial.

```
y = polygen(Integers(7), 'y')
P = y.parent()
p = P([4, 5, 2, 6, 3, 3, 6, 2, 1, 1, 2, 5, 6, 3, 5, 1]); p
```

```
y^15 + 5*y^14 + 3*y^13 + 6*y^12 + 5*y^11 + 2*y^10 + y^9 +
y^8 + 2*y^7 + 6*y^6 + 3*y^5 + 3*y^4 + 6*y^3 + 2*y^2 + 5*y + 4
```

```
p.is_irreducible()
```

```
True
```

```
T.<b> = GF(7^15, modulus=p); T
```

```
Finite Field in b of size 7^15
```

Logarithms in Finite Fields

One useful command we have not described is the `.log()` method for elements of a finite field. Since we now know that the multiplicative group of nonzero elements is cyclic, we can express every element as a power of the generator. The `log` method will return that power.

Usually we will want to use the generator as the base of a lograithm computation in a finite field. However, other bases may be used, wih the understanding that if the base is not a generator, then the logarithm may not exist (i.e. there may not be a solution to the relevant equation).

```
F.<a> = GF(5^4)
a^458
```

```
3*a^3 + 2*a^2 + a + 3
```

```
(3*a^3 + 2*a^2 + a + 3).log(a)
```

```
458
```

```
exponent = (3*a^3 + 2*a^2 + a + 3).log(2*a^3 + 4*a^2 + 4*a)
exponent
```

```
211
```

```
(2*a^3 + 4*a^2 + 4*a)^exponent == 3*a^3 + 2*a^2 + a + 3
```

```
True
```

```
(3*a^3 + 2*a^2 + a + 3).log(a^2 + 4*a + 4)
```

```
Traceback (most recent call last):
...
ValueError: No discrete log of 3*a^3 + 2*a^2 + a + 3 found
to base a^2 + 4*a + 4
```

Since we already know many Sage commands, there is not much else worth introducing before we can work profitably with finite fields. The exercises explore the ways we can examine and exploit the structure of finite fields in Sage.

22.6   Sage Exercises

1.   Create a finite field of order 5^2 and then factor $p(x) = x^{25} - x$ over this field. Comment on what is interesting about this result and why it is not a surprise.

2.   Corollary 22.11 says that the nonzero elements of a finite field are a cyclic group under multiplication. The generator used in Sage is also a generator of this multiplicative group. To see this, create a finite field of order 2^7. Create two lists of the elements of the field: first, use the `.list()` method, then use a list comprehension to generate the proper powers of the generator you specified when you created the field.

   The second list should be the whole field, but will be missing zero. Create the zero element of the field (perhaps by coercing 0 into the field) and `.append()` it to the list of powers. Apply the `sorted()` command to each list and then test the lists for equality.

3.   Subfields of a finite field are completely classified by Theorem 22.7. It is possible to create two finite fields of the correct orders for the supefield/subfield realtionship to hold, and to translate between one and the other. However, in this exercise we will create a subfield of a finite field from scratch. Since the group of nonzero elements in a finite field is cyclic, the nonzero elements of a subfield will form a subgroup of the cyclic group, and necessarily will be cyclic.

   Create a finite field of order 3^6. Theory says there is a subfield of order 3^2, since $2|6$. Determine a generator of multiplicative order 8 for the nonzero elements of this subfield, and construct these 8 elements. Add in the field's zero element to this list. It should be clear that this set of 9 elements is closed under multiplication. Absent our theorems about finite fields and cyclic groups, the closure under addition is not a given. Write a single statement that checks if this set is also closed under addition, by considering all possible sums of elements from the set.

4.   This problem investigates the "separableness" of $\mathbb{Q}(\sqrt{3}, \sqrt{7})$. You can create this number field quickly with the `NumberFieldTower` constructor, along with the polynomials $x^2 - 3$ and $x^2 - 7$. Flatten the tower with the `.absolute_field()` method and use the `.structure()` method to retrieve mappings between the tower and the flattened version. Name the tower `N` and use `a` and `b` as generators. Name the flattened version `L` with `c` as a generator.

   Create a nontrivial ("random") element of `L` using as many powers of `c` as possible (check the degree of `L` to see how many linearly independent powers there are). Request from Sage the minimum polynomial of your random element, thus ensuring the element is a root. Construct the minimum polynomial as a polynomial over `N`, the field tower, and find its factorization. Your factorization should have only linear factors. Each root should be an expression in `a` and `b`, so convert each root into an expression with mathematical notation involving $\sqrt{3}$ and $\sqrt{7}$. Use one of the mappings to verify that one of the roots is indeed the original random element.

   Create a few more random elements, and find a factorization (in `N` or in `L`). For a field to be separable, every element of the field should be a root of *some* separable polynomial. The minimal polynomial is a good polynomial to test. (Why?) Based on the evidence, does it appear that $\mathbb{Q}(\sqrt{3}, \sqrt{7})$ is a separable extension?

5.   Exercise 22.3.21 describes the Frobenius Map, an automorphism of a finite field. If `F` is a finite field in Sage, then `End(F)` will create the automorphism group of `F`, the set of all bijective mappings between the field and itself.

   (a) Work Exercise 22.3.21 to gain an understanding of how and why the Frobenius mapping is a field automorphism. (Do not include any of this in your answer to this

question, but understand that the following will be much easier if you do this problem first.)

(b) For some small, but not trivial, finite fields locate the Frobenius map in the automorphism group. Small might mean $p = 2, 3, 5, 7$ and $3 \leq n \leq 10$, with n prime versus composite.

(c) Once you have located the Frobenius map, describe the other automorphisms. In other words, with a bit of investigation, you should find a description of the automorphisms which will allow you to accurately predict the entire automorphism group for a finite field you have not already explored. (Hint: the automorphism group is a group. What if you "do the operation" between the Frobenius map and itself? Just what is the operation? Try using Sage's multiplicative notation with the elements of the automorphism group.)

(d) What is the "structure" of the automorphism group? What special status does the Frobenius map have in this group?

(e) For any field, the subfield known as the fixed field is an important construction, and will be especially important in the next chapter. Given an automorphism τ of a field E, the subset, $K = \{b \in E \mid \tau(b) = b\}$, can be shown to be a subfield of E. It is known as the **fixed field** of τ in E. For each automorphism of $E = GF(3^6)$ identify the fixed field of the automorphism. Since we understand the structure of subfields of a finite field, it is enough to just determine the order of the fixed field to be able to identify the subfield precisely.

6.   Exercise 22.3.15 suggests that every element of a finite field may be written (expressed) as a sum of squares. This exercise suggests computational experiments which might help you formulate a proof for the exercise.

(a) Construct two small, but not too small, finite fields, one with $p = 2$ and the other with an odd prime. Repeat the following for each field, F.

(b) Choose a "random" element of the field, say $a \in F$. Construct the sets

$$\{x^2 | x \in F\} \qquad\qquad \{a - x^2 | x \in F\}$$

using Sage sets with the $Set()$ constructor. (Be careful: `set()` is a Python command which behaves differently in fundamental ways.)

(c) Examine the size of the two sets and the size of their intersection (`.intersection()`). Try different elements for a, perhaps writing a loop to try *all* possible values. Note that $p = 2$ will behave quite differently.

(d) Suppose you have an element of the intersection. (You can get one with `.an_element()`.) How does this lead to the sum of squares proposed in the exercise?

(e) Can you write a Python function that accepts a finite field whose order is a power of an odd prime and then lists each element as a sum of squares?

Galois Theory

A classic problem of algebra is to find the solutions of a polynomial equation. The solution to the quadratic equation was known in antiquity. Italian mathematicians found general solutions to the general cubic and quartic equations in the sixteenth century; however, attempts to solve the general fifth-degree, or quintic, polynomial were repulsed for the next three hundred years. Certainly, equations such as $x^5 - 1 = 0$ or $x^6 - x^3 - 6 = 0$ could be solved, but no solution like the quadratic formula was found for the general quintic,

$$ax^5 + bx^4 + cx^3 + dx^2 + ex + f = 0.$$

Finally, at the beginning of the nineteenth century, Ruffini and Abel both found quintics that could not be solved with any formula. It was Galois, however, who provided the full explanation by showing which polynomials could and could not be solved by formulas. He discovered the connection between groups and field extensions. Galois theory demonstrates the strong interdependence of group and field theory, and has had far-reaching implications beyond its original purpose.

In this chapter we will prove the Fundamental Theorem of Galois Theory. This result will be used to establish the insolvability of the quintic and to prove the Fundamental Theorem of Algebra.

23.1 Field Automorphisms

Our first task is to establish a link between group theory and field theory by examining automorphisms of fields.

Proposition 23.1. *The set of all automorphisms of a field F is a group under composition of functions.*

PROOF. If σ and τ are automorphisms of F, then so are $\sigma\tau$ and σ^{-1}. The identity is certainly an automorphism; hence, the set of all automorphisms of a field F is indeed a group. $\square$

Proposition 23.2. *Let E be a field extension of F. Then the set of all automorphisms of E that fix F elementwise is a group; that is, the set of all automorphisms $\sigma : E \to E$ such that $\sigma(\alpha) = \alpha$ for all $\alpha \in F$ is a group.*

PROOF. We need only show that the set of automorphisms of E that fix F elementwise is a subgroup of the group of all automorphisms of E. Let σ and τ be two automorphisms of E such that $\sigma(\alpha) = \alpha$ and $\tau(\alpha) = \alpha$ for all $\alpha \in F$. Then $\sigma\tau(\alpha) = \sigma(\alpha) = \alpha$ and $\sigma^{-1}(\alpha) = \alpha$. Since the identity fixes every element of E, the set of automorphisms of E that leave elements of F fixed is a subgroup of the entire group of automorphisms of E. $\square$

Let E be a field extension of F. We will denote the full group of automorphisms of E by $\mathrm{Aut}(E)$. We define the **Galois group** of E over F to be the group of automorphisms of E that fix F elementwise; that is,

$$G(E/F) = \{\sigma \in \mathrm{Aut}(E) : \sigma(\alpha) = \alpha \text{ for all } \alpha \in F\}.$$

If $f(x)$ is a polynomial in $F[x]$ and E is the splitting field of $f(x)$ over F, then we define the Galois group of $f(x)$ to be $G(E/F)$.

Example 23.3. Complex conjugation, defined by $\sigma : a + bi \mapsto a - bi$, is an automorphism of the complex numbers. Since

$$\sigma(a) = \sigma(a + 0i) = a - 0i = a,$$

the automorphism defined by complex conjugation must be in $G(\mathbb{C}/\mathbb{R})$.

Example 23.4. Consider the fields $\mathbb{Q} \subset \mathbb{Q}(\sqrt{5}) \subset \mathbb{Q}(\sqrt{3}, \sqrt{5})$. Then for $a, b \in \mathbb{Q}(\sqrt{5})$,

$$\sigma(a + b\sqrt{3}) = a - b\sqrt{3}$$

is an automorphism of $\mathbb{Q}(\sqrt{3}, \sqrt{5})$ leaving $\mathbb{Q}(\sqrt{5})$ fixed. Similarly,

$$\tau(a + b\sqrt{5}) = a - b\sqrt{5}$$

is an automorphism of $\mathbb{Q}(\sqrt{3}, \sqrt{5})$ leaving $\mathbb{Q}(\sqrt{3})$ fixed. The automorphism $\mu = \sigma\tau$ moves both $\sqrt{3}$ and $\sqrt{5}$. It will soon be clear that $\{\mathrm{id}, \sigma, \tau, \mu\}$ is the Galois group of $\mathbb{Q}(\sqrt{3}, \sqrt{5})$ over $\mathbb{Q}$. The following table shows that this group is isomorphic to $\mathbb{Z}_2 \times \mathbb{Z}_2$.

	id	σ	τ	μ
id	id	σ	τ	μ
σ	σ	id	μ	τ
τ	τ	μ	id	σ
μ	μ	τ	σ	id

We may also regard the field $\mathbb{Q}(\sqrt{3}, \sqrt{5})$ as a vector space over $\mathbb{Q}$ that has basis $\{1, \sqrt{3}, \sqrt{5}, \sqrt{15}\}$. It is no coincidence that $|G(\mathbb{Q}(\sqrt{3}, \sqrt{5})/\mathbb{Q})| = [\mathbb{Q}(\sqrt{3}, \sqrt{5}) : \mathbb{Q}] = 4$.

Proposition 23.5. *Let E be a field extension of F and $f(x)$ be a polynomial in $F[x]$. Then any automorphism in $G(E/F)$ defines a permutation of the roots of $f(x)$ that lie in E.*

PROOF. Let

$$f(x) = a_0 + a_1 x + a_2 x^2 + \cdots + a_n x^n$$

and suppose that $\alpha \in E$ is a zero of $f(x)$. Then for $\sigma \in G(E/F)$,

$$\begin{aligned}
0 &= \sigma(0) \\
&= \sigma(f(\alpha)) \\
&= \sigma(a_0 + a_1\alpha + a_2\alpha^2 + \cdots + a_n\alpha^n) \\
&= a_0 + a_1\sigma(\alpha) + a_2[\sigma(\alpha)]^2 + \cdots + a_n[\sigma(\alpha)]^n;
\end{aligned}$$

therefore, $\sigma(\alpha)$ is also a zero of $f(x)$. $\qquad\square$

Let E be an algebraic extension of a field F. Two elements $\alpha, \beta \in E$ are **conjugate** over F if they have the same minimal polynomial. For example, in the field $\mathbb{Q}(\sqrt{2})$ the elements $\sqrt{2}$ and $-\sqrt{2}$ are conjugate over $\mathbb{Q}$ since they are both roots of the irreducible polynomial $x^2 - 2$.

A converse of the last proposition exists. The proof follows directly from Lemma 21.32.

Proposition 23.6. *If α and β are conjugate over F, there exists an isomorphism $\sigma : F(\alpha) \to F(\beta)$ such that σ is the identity when restricted to F.*

Theorem 23.7. *Let $f(x)$ be a polynomial in $F[x]$ and suppose that E is the splitting field for $f(x)$ over F. If $f(x)$ has no repeated roots, then*

$$|G(E/F)| = [E : F].$$

PROOF. We will use mathematical induction on the degree of $f(x)$. If the degree of $f(x)$ is 0 or 1, then $E = F$ and there is nothing to show. Assume that the result holds for all polynomials of degree k with $0 \leq k < n$. Suppose that the degree of $f(x)$ is n. Let $p(x)$ be an irreducible factor of $f(x)$ of degree r. Since all of the roots of $p(x)$ are in E, we can choose one of these roots, say α, so that $F \subset F(\alpha) \subset E$. Then

$$[E : F(\alpha)] = n/r \quad \text{and} \quad [F(\alpha) : F] = r.$$

If β is any other root of $p(x)$, then $F \subset F(\beta) \subset E$. By Lemma 21.32, there exists a unique isomorphism $\sigma : F(\alpha) \to F(\beta)$ for each such β that fixes F elementwise. Since E is a splitting field of $p(x)$, there are exactly r such isomorphisms. For each of these automorphisms, we can use our induction hypothesis on $[E : F(\alpha)] = n/r < n$ to conclude that

$$|G(E/F(\alpha))| = [E : F(\alpha)].$$

Consequently, there are

$$[E : F] = [E : F(\alpha)][F(\alpha) : F] = n$$

possible automorphisms of E that fix F, or $|G(E/F)| = [E : F]$. $\square$

Corollary 23.8. *Let F be a finite field with a finite extension E such that $[E : F] = k$. Then $G(E/F)$ is cyclic of order k.*

PROOF. Let p be the characteristic of E and F and assume that the orders of E and F are p^m and p^n, respectively. Then $nk = m$. We can also assume that E is the splitting field of $x^{p^m} - x$ over a subfield of order p. Therefore, E must also be the splitting field of $x^{p^m} - x$ over F. Applying Theorem 23.7, we find that $|G(E/F)| = k$.

To prove that $G(E/F)$ is cyclic, we must find a generator for $G(E/F)$. Let $\sigma : E \to E$ be defined by $\sigma(\alpha) = \alpha^{p^n}$. We claim that σ is the element in $G(E/F)$ that we are seeking. We first need to show that σ is in $\text{Aut}(E)$. If α and β are in E,

$$\sigma(\alpha + \beta) = (\alpha + \beta)^{p^n} = \alpha^{p^n} + \beta^{p^n} = \sigma(\alpha) + \sigma(\beta)$$

by Lemma 22.3. Also, it is easy to show that $\sigma(\alpha\beta) = \sigma(\alpha)\sigma(\beta)$. Since σ is a nonzero homomorphism of fields, it must be injective. It must also be onto, since E is a finite field. We know that σ must be in $G(E/F)$, since F is the splitting field of $x^{p^n} - x$ over the base field of order p. This means that σ leaves every element in F fixed. Finally, we must show that the order of σ is k. By Theorem 23.7, we know that

$$\sigma^k(\alpha) = \alpha^{p^{nk}} = \alpha^{p^m} = \alpha$$

is the identity of $G(E/F)$. However, σ^r cannot be the identity for $1 \leq r < k$; otherwise, $x^{p^{nr}} - x$ would have p^m roots, which is impossible. $\square$

Example 23.9. We can now confirm that the Galois group of $\mathbb{Q}(\sqrt{3}, \sqrt{5})$ over $\mathbb{Q}$ in Example 23.4 is indeed isomorphic to $\mathbb{Z}_2 \times \mathbb{Z}_2$. Certainly the group $H = \{\text{id}, \sigma, \tau, \mu\}$ is a subgroup of $G(\mathbb{Q}(\sqrt{3}, \sqrt{5})/\mathbb{Q})$; however, H must be all of $G(\mathbb{Q}(\sqrt{3}, \sqrt{5})/\mathbb{Q})$, since

$$|H| = [\mathbb{Q}(\sqrt{3}, \sqrt{5}) : \mathbb{Q}] = |G(\mathbb{Q}(\sqrt{3}, \sqrt{5})/\mathbb{Q})| = 4.$$

Example 23.10. Let us compute the Galois group of

$$f(x) = x^4 + x^3 + x^2 + x + 1$$

over $\mathbb{Q}$. We know that $f(x)$ is irreducible by Exercise 17.4.20 in Chapter 17. Furthermore, since $(x-1)f(x) = x^5 - 1$, we can use DeMoivre's Theorem to determine that the roots of $f(x)$ are ω^i, where $i = 1, \ldots, 4$ and

$$\omega = \cos(2\pi/5) + i \sin(2\pi/5).$$

Hence, the splitting field of $f(x)$ must be $\mathbb{Q}(\omega)$. We can define automorphisms σ_i of $\mathbb{Q}(\omega)$ by $\sigma_i(\omega) = \omega^i$ for $i = 1, \ldots, 4$. It is easy to check that these are indeed distinct automorphisms in $G(\mathbb{Q}(\omega)/\mathbb{Q})$. Since

$$[\mathbb{Q}(\omega) : \mathbb{Q}] = |G(\mathbb{Q}(\omega)/\mathbb{Q})| = 4,$$

the σ_i's must be all of $G(\mathbb{Q}(\omega)/\mathbb{Q})$. Therefore, $G(\mathbb{Q}(\omega)/\mathbb{Q}) \cong \mathbb{Z}_4$ since ω is a generator for the Galois group.

Separable Extensions

Many of the results that we have just proven depend on the fact that a polynomial $f(x)$ in $F[x]$ has no repeated roots in its splitting field. It is evident that we need to know exactly when a polynomial factors into distinct linear factors in its splitting field. Let E be the splitting field of a polynomial $f(x)$ in $F[x]$. Suppose that $f(x)$ factors over E as

$$f(x) = (x - \alpha_1)^{n_1}(x - \alpha_2)^{n_2} \cdots (x - \alpha_r)^{n_r} = \prod_{i=1}^{r}(x - \alpha_i)^{n_i}.$$

We define the **multiplicity** of a root α_i of $f(x)$ to be n_i. A root with multiplicity 1 is called a **simple root**. Recall that a polynomial $f(x) \in F[x]$ of degree n is **separable** if it has n distinct roots in its splitting field E. Equivalently, $f(x)$ is separable if it factors into distinct linear factors over $E[x]$. An extension E of F is a **separable extension** of F if every element in E is the root of a separable polynomial in $F[x]$. Also recall that $f(x)$ is separable if and only if $\gcd(f(x), f'(x)) = 1$ (Lemma 22.5).

Proposition 23.11. *Let $f(x)$ be an irreducible polynomial over F. If the characteristic of F is 0, then $f(x)$ is separable. If the characteristic of F is p and $f(x) \neq g(x^p)$ for some $g(x)$ in $F[x]$, then $f(x)$ is also separable.*

PROOF.   First assume that $\operatorname{char} F = 0$. Since $\deg f'(x) < \deg f(x)$ and $f(x)$ is irreducible, the only way $\gcd(f(x), f'(x)) \neq 1$ is if $f'(x)$ is the zero polynomial; however, this is impossible in a field of characteristic zero. If $\operatorname{char} F = p$, then $f'(x)$ can be the zero polynomial if every coefficient of $f'(x)$ is a multiple of p. This can happen only if we have a polynomial of the form $f(x) = a_0 + a_1 x^p + a_2 x^{2p} + \cdots + a_n x^{np}$.   □

Certainly extensions of a field F of the form $F(\alpha)$ are some of the easiest to study and understand. Given a field extension E of F, the obvious question to ask is when it is possible to find an element $\alpha \in E$ such that $E = F(\alpha)$. In this case, α is called a **primitive element**. We already know that primitive elements exist for certain extensions. For example,

$$\mathbb{Q}(\sqrt{3}, \sqrt{5}) = \mathbb{Q}(\sqrt{3} + \sqrt{5})$$

and

$$\mathbb{Q}(\sqrt[3]{5}, \sqrt{5}\,i) = \mathbb{Q}(\sqrt[6]{5}\,i).$$

Corollary 22.12 tells us that there exists a primitive element for any finite extension of a finite field. The next theorem tells us that we can often find a primitive element.

Theorem 23.12 Primitive Element Theorem. *Let E be a finite separable extension of a field F. Then there exists an $\alpha \in E$ such that $E = F(\alpha)$.*

PROOF. We already know that there is no problem if F is a finite field. Suppose that E is a finite extension of an infinite field. We will prove the result for $F(\alpha, \beta)$. The general case easily follows when we use mathematical induction. Let $f(x)$ and $g(x)$ be the minimal polynomials of α and β, respectively. Let K be the field in which both $f(x)$ and $g(x)$ split. Suppose that $f(x)$ has zeros $\alpha = \alpha_1, \ldots, \alpha_n$ in K and $g(x)$ has zeros $\beta = \beta_1, \ldots, \beta_m$ in K. All of these zeros have multiplicity 1, since E is separable over F. Since F is infinite, we can find an a in F such that

$$a \neq \frac{\alpha_i - \alpha}{\beta - \beta_j}$$

for all i and j with $j \neq 1$. Therefore, $a(\beta - \beta_j) \neq \alpha_i - \alpha$. Let $\gamma = \alpha + a\beta$. Then

$$\gamma = \alpha + a\beta \neq \alpha_i + a\beta_j;$$

hence, $\gamma - a\beta_j \neq \alpha_i$ for all i, j with $j \neq 1$. Define $h(x) \in F(\gamma)[x]$ by $h(x) = f(\gamma - ax)$. Then $h(\beta) = f(\alpha) = 0$. However, $h(\beta_j) \neq 0$ for $j \neq 1$. Hence, $h(x)$ and $g(x)$ have a single common factor in $F(\gamma)[x]$; that is, the minimal polynomial of β over $F(\gamma)$ must be linear, since β is the only zero common to both $g(x)$ and $h(x)$. So $\beta \in F(\gamma)$ and $\alpha = \gamma - a\beta$ is in $F(\gamma)$. Hence, $F(\alpha, \beta) = F(\gamma)$. $\qquad\square$

23.2 The Fundamental Theorem

The goal of this section is to prove the Fundamental Theorem of Galois Theory. This theorem explains the connection between the subgroups of $G(E/F)$ and the intermediate fields between E and F.

Proposition 23.13. *Let $\{\sigma_i : i \in I\}$ be a collection of automorphisms of a field F. Then*

$$F_{\{\sigma_i\}} = \{a \in F : \sigma_i(a) = a \text{ for all } \sigma_i\}$$

is a subfield of F.

PROOF. Let $\sigma_i(a) = a$ and $\sigma_i(b) = b$. Then

$$\sigma_i(a \pm b) = \sigma_i(a) \pm \sigma_i(b) = a \pm b$$

and

$$\sigma_i(ab) = \sigma_i(a)\sigma_i(b) = ab.$$

If $a \neq 0$, then $\sigma_i(a^{-1}) = [\sigma_i(a)]^{-1} = a^{-1}$. Finally, $\sigma_i(0) = 0$ and $\sigma_i(1) = 1$ since σ_i is an automorphism. $\qquad\square$

Corollary 23.14. *Let F be a field and let G be a subgroup of $\mathrm{Aut}(F)$. Then*

$$F_G = \{\alpha \in F : \sigma(\alpha) = \alpha \text{ for all } \sigma \in G\}$$

is a subfield of F.

The subfield $F_{\{\sigma_i\}}$ of F is called the **fixed field** of $\{\sigma_i\}$. The field fixed by a subgroup G of $\mathrm{Aut}(F)$ will be denoted by F_G.

Example 23.15. Let $\sigma : \mathbb{Q}(\sqrt{3}, \sqrt{5}) \to \mathbb{Q}(\sqrt{3}, \sqrt{5})$ be the automorphism that maps $\sqrt{3}$ to $-\sqrt{3}$. Then $\mathbb{Q}(\sqrt{5})$ is the subfield of $\mathbb{Q}(\sqrt{3}, \sqrt{5})$ left fixed by σ.

Proposition 23.16. *Let E be a splitting field over F of a separable polynomial. Then $E_{G(E/F)} = F$.*

PROOF.   Let $G = G(E/F)$. Clearly, $F \subset E_G \subset E$. Also, E must be a splitting field of E_G and $G(E/F) = G(E/E_G)$. By Theorem 23.7,

$$|G| = [E : E_G] = [E : F].$$

Therefore, $[E_G : F] = 1$. Consequently, $E_G = F$. $\square$

A large number of mathematicians first learned Galois theory from Emil Artin's monograph on the subject [1]. The very clever proof of the following lemma is due to Artin.

Lemma 23.17. *Let G be a finite group of automorphisms of E and let $F = E_G$. Then $[E : F] \leq |G|$.*

PROOF.   Let $|G| = n$. We must show that any set of $n + 1$ elements $\alpha_1, \ldots, \alpha_{n+1}$ in E is linearly dependent over F; that is, we need to find elements $a_i \in F$, not all zero, such that

$$a_1\alpha_1 + a_2\alpha_2 + \cdots + a_{n+1}\alpha_{n+1} = 0.$$

Suppose that $\sigma_1 = \mathrm{id}, \sigma_2, \ldots, \sigma_n$ are the automorphisms in G. The homogeneous system of linear equations

$$\sigma_1(\alpha_1)x_1 + \sigma_1(\alpha_2)x_2 + \cdots + \sigma_1(\alpha_{n+1})x_{n+1} = 0$$
$$\sigma_2(\alpha_1)x_1 + \sigma_2(\alpha_2)x_2 + \cdots + \sigma_2(\alpha_{n+1})x_{n+1} = 0$$
$$\vdots$$
$$\sigma_n(\alpha_1)x_1 + \sigma_n(\alpha_2)x_2 + \cdots + \sigma_n(\alpha_{n+1})x_{n+1} = 0$$

has more unknowns than equations. From linear algebra we know that this system has a nontrivial solution, say $x_i = a_i$ for $i = 1, 2, \ldots, n + 1$. Since σ_1 is the identity, the first equation translates to

$$a_1\alpha_1 + a_2\alpha_2 + \cdots + a_{n+1}\alpha_{n+1} = 0.$$

The problem is that some of the a_i's may be in E but not in F. We must show that this is impossible.

Suppose that at least one of the a_i's is in E but not in F. By rearranging the α_i's we may assume that a_1 is nonzero. Since any nonzero multiple of a solution is also a solution, we can also assume that $a_1 = 1$. Of all possible solutions fitting this description, we choose the one with the smallest number of nonzero terms. Again, by rearranging $\alpha_2, \ldots, \alpha_{n+1}$ if necessary, we can assume that a_2 is in E but not in F. Since F is the subfield of E that is fixed elementwise by G, there exists a σ_i in G such that $\sigma_i(a_2) \neq a_2$. Applying σ_i to each equation in the system, we end up with the same homogeneous system, since G is a group. Therefore, $x_1 = \sigma_i(a_1) = 1$, $x_2 = \sigma_i(a_2)$, $\ldots$, $x_{n+1} = \sigma_i(a_{n+1})$ is also a solution of the original system. We know that a linear combination of two solutions of a homogeneous system is also a solution; consequently,

$$x_1 = 1 - 1 = 0$$
$$x_2 = a_2 - \sigma_i(a_2)$$
$$\vdots$$
$$x_{n+1} = a_{n+1} - \sigma_i(a_{n+1})$$

must be another solution of the system. This is a nontrivial solution because $\sigma_i(a_2) \neq a_2$, and has fewer nonzero entries than our original solution. This is a contradiction, since the number of nonzero solutions to our original solution was assumed to be minimal. We can therefore conclude that $a_1, \ldots, a_{n+1} \in F$. $\qquad \square$

Let E be an algebraic extension of F. If every irreducible polynomial in $F[x]$ with a root in E has all of its roots in E, then E is called a ***normal extension*** of F; that is, every irreducible polynomial in $F[x]$ containing a root in E is the product of linear factors in $E[x]$.

Theorem 23.18. *Let E be a field extension of F. Then the following statements are equivalent.*

1. *E is a finite, normal, separable extension of F.*

2. *E is a splitting field over F of a separable polynomial.*

3. *$F = E_G$ for some finite group G of automorphisms of E.*

PROOF. $(1) \Rightarrow (2)$. Let E be a finite, normal, separable extension of F. By the Primitive Element Theorem, we can find an α in E such that $E = F(\alpha)$. Let $f(x)$ be the minimal polynomial of α over F. The field E must contain all of the roots of $f(x)$ since it is a normal extension F; hence, E is a splitting field for $f(x)$.

$(2) \Rightarrow (3)$. Let E be the splitting field over F of a separable polynomial. By Proposition 23.16, $E_{G(E/F)} = F$. Since $|G(E/F)| = [E : F]$, this is a finite group.

$(3) \Rightarrow (1)$. Let $F = E_G$ for some finite group of automorphisms G of E. Since $[E : F] \leq |G|$, E is a finite extension of F. To show that E is a finite, normal extension of F, let $f(x) \in F[x]$ be an irreducible monic polynomial that has a root α in E. We must show that $f(x)$ is the product of distinct linear factors in $E[x]$. By Proposition 23.5, automorphisms in G permute the roots of $f(x)$ lying in E. Hence, if we let G act on α, we can obtain distinct roots $\alpha_1 = \alpha, \alpha_2, \ldots, \alpha_n$ in E. Let $g(x) = \prod_{i=1}^{n}(x - \alpha_i)$. Then $g(x)$ is separable over F and $g(\alpha) = 0$. Any automorphism σ in G permutes the factors of $g(x)$ since it permutes these roots; hence, when σ acts on $g(x)$, it must fix the coefficients of $g(x)$. Therefore, the coefficients of $g(x)$ must be in F. Since $\deg g(x) \leq \deg f(x)$ and $f(x)$ is the minimal polynomial of α, $f(x) = g(x)$. $\qquad \square$

Corollary 23.19. *Let K be a field extension of F such that $F = K_G$ for some finite group of automorphisms G of K. Then $G = G(K/F)$.*

PROOF. Since $F = K_G$, G is a subgroup of $G(K/F)$. Hence,

$$[K : F] \leq |G| \leq |G(K/F)| = [K : F].$$

It follows that $G = G(K/F)$, since they must have the same order. $\qquad \square$

Before we determine the exact correspondence between field extensions and automorphisms of fields, let us return to a familiar example.

Example 23.20. In Example 23.4 we examined the automorphisms of $\mathbb{Q}(\sqrt{3}, \sqrt{5})$ fixing $\mathbb{Q}$. Figure 23.21 compares the lattice of field extensions of $\mathbb{Q}$ with the lattice of subgroups of $G(\mathbb{Q}(\sqrt{3}, \sqrt{5})/\mathbb{Q})$. The Fundamental Theorem of Galois Theory tells us what the relationship is between the two lattices.
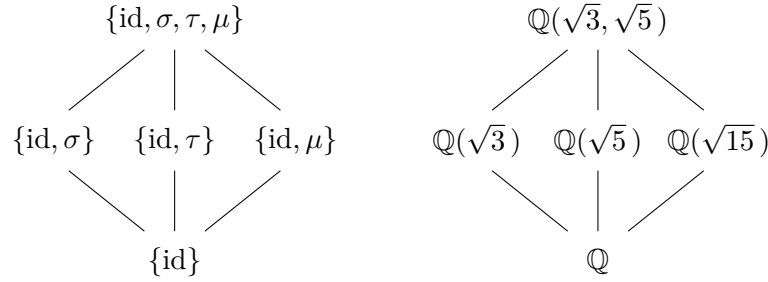
Figure 23.21: $G(\mathbb{Q}(\sqrt{3}, \sqrt{5})/\mathbb{Q})$

We are now ready to state and prove the Fundamental Theorem of Galois Theory.

Theorem 23.22 Fundamental Theorem of Galois Theory. *Let F be a finite field or a field of characteristic zero. If E is a finite normal extension of F with Galois group $G(E/F)$, then the following statements are true.*

1. *The map $K \mapsto G(E/K)$ is a bijection of subfields K of E containing F with the subgroups of $G(E/F)$.*

2. *If $F \subset K \subset E$, then*

$$[E : K] = |G(E/K)| \ and \ [K : F] = [G(E/F) : G(E/K)].$$

3. *$F \subset K \subset L \subset E$ if and only if $\{id\} \subset G(E/L) \subset G(E/K) \subset G(E/F)$.*

4. *K is a normal extension of F if and only if $G(E/K)$ is a normal subgroup of $G(E/F)$. In this case*

$$G(K/F) \cong G(E/F)/G(E/K).$$

PROOF.   (1) Suppose that $G(E/K) = G(E/L) = G$. Both K and L are fixed fields of G; hence, $K = L$ and the map defined by $K \mapsto G(E/K)$ is one-to-one. To show that the map is onto, let G be a subgroup of $G(E/F)$ and K be the field fixed by G. Then $F \subset K \subset E$; consequently, E is a normal extension of K. Thus, $G(E/K) = G$ and the map $K \mapsto G(E/K)$ is a bijection.

(2) By Theorem 23.7, $|G(E/K)| = [E : K]$; therefore,

$$|G(E/F)| = [G(E/F) : G(E/K)] \cdot |G(E/K)| = [E : F] = [E : K][K : F].$$

Thus, $[K : F] = [G(E/F) : G(E/K)]$.

Statement (3) is illustrated in Figure 23.23. We leave the proof of this property as an exercise.

(4) This part takes a little more work. Let K be a normal extension of F. If σ is in $G(E/F)$ and τ is in $G(E/K)$, we need to show that $\sigma^{-1}\tau\sigma$ is in $G(E/K)$; that is, we need to show that $\sigma^{-1}\tau\sigma(\alpha) = \alpha$ for all $\alpha \in K$. Suppose that $f(x)$ is the minimal polynomial of α over F. Then $\sigma(\alpha)$ is also a root of $f(x)$ lying in K, since K is a normal extension of F. Hence, $\tau(\sigma(\alpha)) = \sigma(\alpha)$ or $\sigma^{-1}\tau\sigma(\alpha) = \alpha$.

Conversely, let $G(E/K)$ be a normal subgroup of $G(E/F)$. We need to show that $F = K_{G(K/F)}$. Let $\tau \in G(E/K)$. For all $\sigma \in G(E/F)$ there exists a $\overline{\tau} \in G(E/K)$ such that $\tau\sigma = \sigma\overline{\tau}$. Consequently, for all $\alpha \in K$

$$\tau(\sigma(\alpha)) = \sigma(\overline{\tau}(\alpha)) = \sigma(\alpha);$$

hence, $\sigma(\alpha)$ must be in the fixed field of $G(E/K)$. Let $\overline{\sigma}$ be the restriction of σ to K. Then $\overline{\sigma}$ is an automorphism of K fixing F, since $\sigma(\alpha) \in K$ for all $\alpha \in K$; hence, $\overline{\sigma} \in G(K/F)$. Next, we will show that the fixed field of $G(K/F)$ is F. Let β be an element in K that is fixed by all automorphisms in $G(K/F)$. In particular, $\overline{\sigma}(\beta) = \beta$ for all $\sigma \in G(E/F)$. Therefore, β belongs to the fixed field F of $G(E/F)$.

Finally, we must show that when K is a normal extension of F,

$$G(K/F) \cong G(E/F)/G(E/K).$$

For $\sigma \in G(E/F)$, let σ_K be the automorphism of K obtained by restricting σ to K. Since K is a normal extension, the argument in the preceding paragraph shows that $\sigma_K \in G(K/F)$. Consequently, we have a map $\phi : G(E/F) \to G(K/F)$ defined by $\sigma \mapsto \sigma_K$. This map is a group homomorphism since

$$\phi(\sigma\tau) = (\sigma\tau)_K = \sigma_K\tau_K = \phi(\sigma)\phi(\tau).$$

The kernel of ϕ is $G(E/K)$. By (2),

$$|G(E/F)|/|G(E/K)| = [K : F] = |G(K/F)|.$$

Hence, the image of ϕ is $G(K/F)$ and ϕ is onto. Applying the First Isomorphism Theorem, we have

$$G(K/F) \cong G(E/F)/G(E/K). \qquad \square$$

$$
\begin{array}{ccc}
E & \longrightarrow & \{\text{id}\} \\
| & & | \\
L & \longrightarrow & G(E/L) \\
| & & | \\
K & \longrightarrow & G(E/K) \\
| & & | \\
F & \longrightarrow & G(E/F)
\end{array}
$$
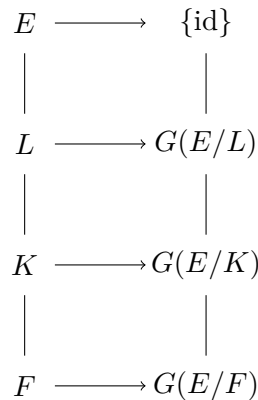
Figure 23.23: Subgroups of $G(E/F)$ and subfields of E

Example 23.24. In this example we will illustrate the Fundamental Theorem of Galois Theory by determining the lattice of subgroups of the Galois group of $f(x) = x^4 - 2$. We will compare this lattice to the lattice of field extensions of $\mathbb{Q}$ that are contained in the splitting field of $x^4 - 2$. The splitting field of $f(x)$ is $\mathbb{Q}(\sqrt[4]{2}, i)$. To see this, notice that $f(x)$ factors as $(x^2 + \sqrt{2})(x^2 - \sqrt{2})$; hence, the roots of $f(x)$ are $\pm\sqrt[4]{2}$ and $\pm\sqrt[4]{2}\,i$. We first adjoin the root $\sqrt[4]{2}$ to $\mathbb{Q}$ and then adjoin the root i of $x^2 + 1$ to $\mathbb{Q}(\sqrt[4]{2})$. The splitting field of $f(x)$ is then $\mathbb{Q}(\sqrt[4]{2})(i) = \mathbb{Q}(\sqrt[4]{2}, i)$.

Since $[\mathbb{Q}(\sqrt[4]{2}) : \mathbb{Q}] = 4$ and i is not in $\mathbb{Q}(\sqrt[4]{2})$, it must be the case that $[\mathbb{Q}(\sqrt[4]{2}, i) : \mathbb{Q}(\sqrt[4]{2})] = 2$. Hence, $[\mathbb{Q}(\sqrt[4]{2}, i) : \mathbb{Q}] = 8$. The set

$$\{1, \sqrt[4]{2}, (\sqrt[4]{2})^2, (\sqrt[4]{2})^3, i, i\sqrt[4]{2}, i(\sqrt[4]{2})^2, i(\sqrt[4]{2})^3\}$$

is a basis of $\mathbb{Q}(\sqrt[4]{2}, i)$ over $\mathbb{Q}$. The lattice of field extensions of $\mathbb{Q}$ contained in $\mathbb{Q}(\sqrt[4]{2}, i)$ is illustrated in Figure 23.25(a).

The Galois group G of $f(x)$ must be of order 8. Let σ be the automorphism defined by $\sigma(\sqrt[4]{2}) = i\sqrt[4]{2}$ and $\sigma(i) = i$, and τ be the automorphism defined by complex conjugation;

that is, $\tau(i) = -i$. Then G has an element of order 4 and an element of order 2. It is easy to verify by direct computation that the elements of G are $\{\text{id}, \sigma, \sigma^2, \sigma^3, \tau, \sigma\tau, \sigma^2\tau, \sigma^3\tau\}$ and that the relations $\tau^2 = \text{id}$, $\sigma^4 = \text{id}$, and $\tau\sigma\tau = \sigma^{-1}$ are satisfied; hence, G must be isomorphic to D_4. The lattice of subgroups of G is illustrated in Figure 23.25(b).
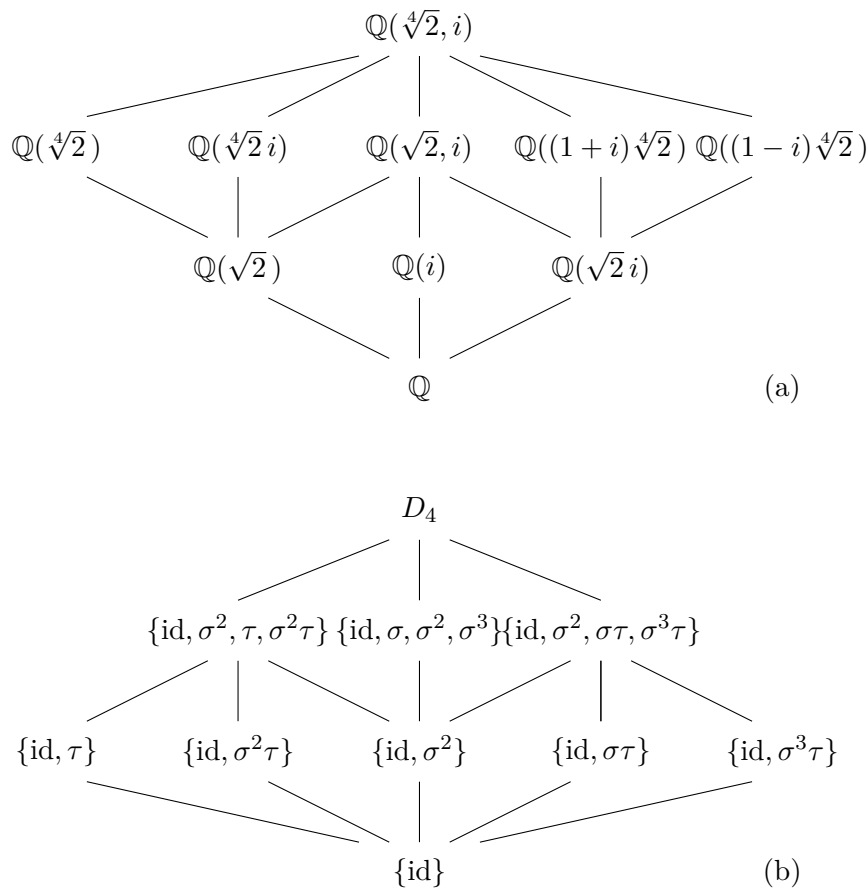
Figure 23.25: Galois group of $x^4 - 2$

<div align="center">

─────────────   ***Historical Note***   ─────────────

</div>

Solutions for the cubic and quartic equations were discovered in the 1500s. Attempts to find solutions for the quintic equations puzzled some of history's best mathematicians. In 1798, P. Ruffini submitted a paper that claimed no such solution could be found; however, the paper was not well received. In 1826, Niels Henrik Abel (1802–1829) finally offered the first correct proof that quintics are not always solvable by radicals.

Abel inspired the work of Évariste Galois. Born in 1811, Galois began to display extraordinary mathematical talent at the age of 14. He applied for entrance to the École Polytechnique several times; however, he had great difficulty meeting the formal entrance requirements, and the examiners failed to recognize his mathematical genius. He was finally accepted at the École Normale in 1829.

Galois worked to develop a theory of solvability for polynomials. In 1829, at the age of 17, Galois presented two papers on the solution of algebraic equations to the Académie des Sciences de Paris. These papers were sent to Cauchy, who subsequently lost them. A third paper was submitted to Fourier, who died before he could read the paper. Another paper

was presented, but was not published until 1846.

Galois' democratic sympathies led him into the Revolution of 1830. He was expelled from school and sent to prison for his part in the turmoil. After his release in 1832, he was drawn into a duel possibly over a love affair. Certain that he would be killed, he spent the evening before his death outlining his work and his basic ideas for research in a long letter to his friend Chevalier. He was indeed dead the next day, at the age of 20.

23.3   Applications

Solvability by Radicals

Throughout this section we shall assume that all fields have characteristic zero to ensure that irreducible polynomials do not have multiple roots. The immediate goal of this section is to determine when the roots of a polynomial $f(x)$ can be computed with a finite number of operations on the coefficients of $f(x)$. The allowable operations are addition, subtraction, multiplication, division, and the extraction of nth roots. Certainly the solution to the quadratic equation, $ax^2 + bx + c = 0$, illustrates this process:

$$x = \frac{-b \pm \sqrt{b^2 - 4ac}}{2a}.$$

The only one of these operations that might demand a larger field is the taking of nth roots. We are led to the following definition.

An extension field E of a field F is an **extension by radicals** if there exists a chain of subfields

$$F = F_0 \subset F_1 \subset F_2 \subset \cdots \subset F_r = E$$

such for $i = 1, 2, \ldots, r$, we have $F_i = F_{i-1}(\alpha_i)$ and $\alpha_i^{n_i} \in F_{i-1}$ for some positive integer n_i. A polynomial $f(x)$ is **solvable by radicals** over F if the splitting field K of $f(x)$ over F is contained in an extension of F by radicals. Our goal is to arrive at criteria that will tell us whether or not a polynomial $f(x)$ is solvable by radicals by examining the Galois group $f(x)$.

The easiest polynomial to solve by radicals is one of the form $x^n - a$. As we discussed in Chapter 4, the roots of $x^n - 1$ are called the **nth roots of unity**. These roots are a finite subgroup of the splitting field of $x^n - 1$. By Corollary 22.11, the nth roots of unity form a cyclic group. Any generator of this group is called a **primitive nth root of unity**.

Example 23.26. The polynomial $x^n - 1$ is solvable by radicals over $\mathbb{Q}$. The roots of this polynomial are $1, \omega, \omega^2, \ldots, \omega^{n-1}$, where

$$\omega = \cos\left(\frac{2\pi}{n}\right) + i \sin\left(\frac{2\pi}{n}\right).$$

The splitting field of $x^n - 1$ over $\mathbb{Q}$ is $\mathbb{Q}(\omega)$.

We shall prove that a polynomial is solvable by radicals if its Galois group is solvable. Recall that a subnormal series of a group G is a finite sequence of subgroups

$$G = H_n \supset H_{n-1} \supset \cdots \supset H_1 \supset H_0 = \{e\},$$

where H_i is normal in H_{i+1}. A group G is solvable if it has a subnormal series $\{H_i\}$ such that all of the factor groups H_{i+1}/H_i are abelian. For example, if we examine the series $\{\mathrm{id}\} \subset A_3 \subset S_3$, we see that S_3 is solvable. On the other hand, S_5 is not solvable, by Theorem 10.11.

Lemma 23.27. *Let F be a field of characteristic zero and E be the splitting field of $x^n - a$ over F with $a \in F$. Then $G(E/F)$ is a solvable group.*

PROOF.   The roots of $x^n - a$ are $\sqrt[n]{a}, \omega \sqrt[n]{a}, \ldots, \omega^{n-1} \sqrt[n]{a}$, where ω is a primitive nth root of unity. Suppose that F contains all of its nth roots of unity. If ζ is one of the roots of $x^n - a$, then distinct roots of $x^n - a$ are $\zeta, \omega\zeta, \ldots, \omega^{n-1}\zeta$, and $E = F(\zeta)$. Since $G(E/F)$ permutes the roots $x^n - a$, the elements in $G(E/F)$ must be determined by their action on these roots. Let σ and τ be in $G(E/F)$ and suppose that $\sigma(\zeta) = \omega^i \zeta$ and $\tau(\zeta) = \omega^j \zeta$. If F contains the roots of unity, then

$$\sigma\tau(\zeta) = \sigma(\omega^j \zeta) = \omega^j \sigma(\zeta) = \omega^{i+j}\zeta = \omega^i \tau(\zeta) = \tau(\omega^i \zeta) = \tau\sigma(\zeta).$$

Therefore, $\sigma\tau = \tau\sigma$ and $G(E/F)$ is abelian, and $G(E/F)$ must be solvable.

   Now suppose that F does not contain a primitive nth root of unity. Let ω be a generator of the cyclic group of the nth roots of unity. Let α be a zero of $x^n - a$. Since α and $\omega\alpha$ are both in the splitting field of $x^n - a$, $\omega = (\omega\alpha)/\alpha$ is also in E. Let $K = F(\omega)$. Then $F \subset K \subset E$. Since K is the splitting field of $x^n - 1$, K is a normal extension of F. Therefore, any automorphism σ in $G(F(\omega)/F)$ is determined by $\sigma(\omega)$. It must be the case that $\sigma(\omega) = \omega^i$ for some integer i since all of the zeros of $x^n - 1$ are powers of ω. If $\tau(\omega) = \omega^j$ is in $G(F(\omega)/F)$, then

$$\sigma\tau(\omega) = \sigma(\omega^j) = [\sigma(\omega)]^j = \omega^{ij} = [\tau(\omega)]^i = \tau(\omega^i) = \tau\sigma(\omega).$$

Therefore, $G(F(\omega)/F)$ is abelian. By the Fundamental Theorem of Galois Theory the series

$$\{\mathrm{id}\} \subset G(E/F(\omega)) \subset G(E/F)$$

is a normal series. By our previous argument, $G(E/F(\omega))$ is abelian. Since

$$G(E/F)/G(E/F(\omega)) \cong G(F(\omega)/F)$$

is also abelian, $G(E/F)$ is solvable. $\square$

Lemma 23.28. *Let F be a field of characteristic zero and let*

$$F = F_0 \subset F_1 \subset F_2 \subset \cdots \subset F_r = E$$

a radical extension of F. Then there exists a normal radical extension

$$F = K_0 \subset K_1 \subset K_2 \subset \cdots \subset K_r = K$$

such that K that contains E and K_i is a normal extension of K_{i-1}.

PROOF.   Since E is a radical extension of F, there exists a chain of subfields

$$F = F_0 \subset F_1 \subset F_2 \subset \cdots \subset F_r = E$$

such for $i = 1, 2, \ldots, r$, we have $F_i = F_{i-1}(\alpha_i)$ and $\alpha_i^{n_i} \in F_{i-1}$ for some positive integer n_i. We will build a normal radical extension of F,

$$F = K_0 \subset K_1 \subset K_2 \subset \cdots \subset K_r = K$$

such that $K \supseteq E$. Define K_1 for be the splitting field of $x^{n_1} - \alpha_1^{n_1}$. The roots of this polynomial are $\alpha_1, \alpha_1\omega, \alpha_1\omega^2, \ldots, \alpha_1\omega^{n_1 - 1}$, where ω is a primitive n_1th root of unity. If F contains all of its n_1 roots of unity, then $K_1 = F(\alpha_!)$. On the other hand, suppose that F does not contain a primitive n_1th root of unity. If β is a root of $x^{n_1} - \alpha_1^{n_1}$, then all of

the roots of $x^{n_1} - \alpha_1^{n_1}$ must be $\beta, \omega\beta, \ldots, \omega^{n_1-1}$, where ω is a primitive n_1th root of unity. In this case, $K_1 = F(\omega\beta)$. Thus, K_1 is a normal radical extension of F containing F_1. Continuing in this manner, we obtain

$$F = K_0 \subset K_1 \subset K_2 \subset \cdots \subset K_r = K$$

such that K_i is a normal extension of K_{i-1} and $K_i \supseteq F_i$ for $i = 1, 2, \ldots, r$. $\qquad \square$

We will now prove the main theorem about solvability by radicals.

Theorem 23.29. *Let $f(x)$ be in $F[x]$, where $\operatorname{char} F = 0$. If $f(x)$ is solvable by radicals, then the Galois group of $f(x)$ over F is solvable.*

PROOF. Since $f(x)$ is solvable by radicals there exists an extension E of F by radicals $F = F_0 \subset F_1 \subset \cdots \subset F_n = E$. By Lemma 23.28, we can assume that E is a splitting field $f(x)$ and F_i is normal over F_{i-1}. By the Fundamental Theorem of Galois Theory, $G(E/F_i)$ is a normal subgroup of $G(E/F_{i-1})$. Therefore, we have a subnormal series of subgroups of $G(E/F)$:

$$\{\operatorname{id}\} \subset G(E/F_{n-1}) \subset \cdots \subset G(E/F_1) \subset G(E/F).$$

Again by the Fundamental Theorem of Galois Theory, we know that

$$G(E/F_{i-1})/G(E/F_i) \cong G(F_i/F_{i-1}).$$

By Lemma 23.27, $G(F_i/F_{i-1})$ is solvable; hence, $G(E/F)$ is also solvable. $\qquad \square$

The converse of Theorem 23.29 is also true. For a proof, see any of the references at the end of this chapter.

Insolvability of the Quintic

We are now in a position to find a fifth-degree polynomial that is not solvable by radicals. We merely need to find a polynomial whose Galois group is S_5. We begin by proving a lemma.

Lemma 23.30. *If p is prime, then any subgroup of S_p that contains a transposition and a cycle of length p must be all of S_p.*

PROOF. Let G be a subgroup of S_p that contains a transposition σ and τ a cycle of length p. We may assume that $\sigma = (12)$. The order of τ is p and τ^n must be a cycle of length p for $1 \le n < p$. Therefore, we may assume that $\mu = \tau^n = (12i_3 \ldots i_p)$ for some n, where $1 \le n < p$ (see Exercise 5.3.13 in Chapter 5). Noting that $(12)(12i_3 \ldots i_p) = (2i_3 \ldots i_p)$ and $(2i_3 \ldots i_p)^k(12)(2i_3 \ldots i_p)^{-k} = (1i_k)$, we can obtain all the transpositions of the form $(1n)$ for $1 \le n < p$. However, these transpositions generate all transpositions in S_p, since $(1j)(1i)(1j) = (ij)$. The transpositions generate S_p. $\qquad \square$
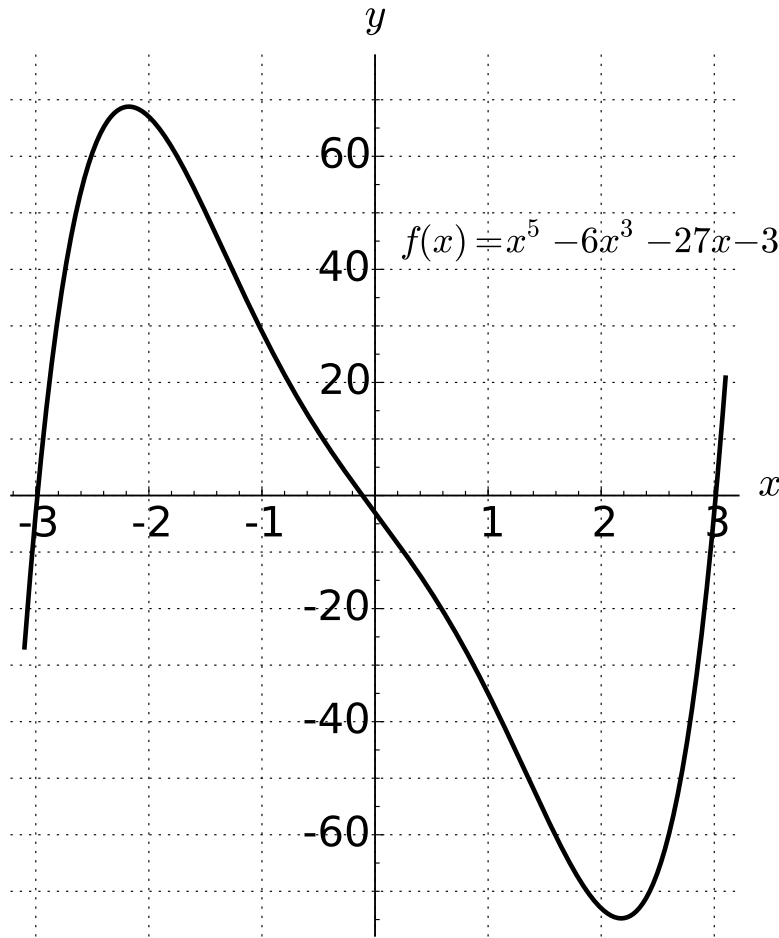
Figure 23.31: The graph of $f(x) = x^5 - 6x^3 - 27x - 3$

Example 23.32. We will show that $f(x) = x^5 - 6x^3 - 27x - 3 \in \mathbb{Q}[x]$ is not solvable. We claim that the Galois group of $f(x)$ over $\mathbb{Q}$ is S_5. By Eisenstein's Criterion, $f(x)$ is irreducible and, therefore, must be separable. The derivative of $f(x)$ is $f'(x) = 5x^4 - 18x^2 - 27$; hence, setting $f'(x) = 0$ and solving, we find that the only real roots of $f'(x)$ are

$$x = \pm\sqrt{\frac{6\sqrt{6} + 9}{5}}.$$

Therefore, $f(x)$ can have at most one maximum and one minimum. It is easy to show that $f(x)$ changes sign between -3 and -2, between -2 and 0, and once again between 0 and 4 (Figure 23.31). Therefore, $f(x)$ has exactly three distinct real roots. The remaining two roots of $f(x)$ must be complex conjugates. Let K be the splitting field of $f(x)$. Since $f(x)$ has five distinct roots in K and every automorphism of K fixing $\mathbb{Q}$ is determined by the way it permutes the roots of $f(x)$, we know that $G(K/\mathbb{Q})$ is a subgroup of S_5. Since f is irreducible, there is an element in $\sigma \in G(K/\mathbb{Q})$ such that $\sigma(a) = b$ for two roots a and b of $f(x)$. The automorphism of $\mathbb{C}$ that takes $a + bi \mapsto a - bi$ leaves the real roots fixed and interchanges the complex roots; consequently, $G(K/\mathbb{Q})$ contains a transpostion. If α is one of the real roots of $f(x)$, then $[\mathbb{Q}(\alpha) : \mathbb{Q}] = 5$ by Exercise 21.4.28. Since $\mathbb{Q}(\alpha)$ is a subfield of K, it must be the case the $[K : \mathbb{Q}]$ is divisible by 5. Since $[K : \mathbb{Q}] = |G(K/\mathbb{Q})|$ and $G(K/\mathbb{Q}) \subset S_5$, we know that $G(K/\mathbb{Q})$ contains a cycle of length 5. By Lemma 23.30, S_5 is generated by a transposition and an element of order 5; therefore, $G(K/\mathbb{Q})$ must be all of

S_5. By Theorem 10.11, S_5 is not solvable. Consequently, $f(x)$ cannot be solved by radicals.

The Fundamental Theorem of Algebra

It seems fitting that the last theorem that we will state and prove is the Fundamental Theorem of Algebra. This theorem was first proven by Gauss in his doctoral thesis. Prior to Gauss's proof, mathematicians suspected that there might exist polynomials over the real and complex numbers having no solutions. The Fundamental Theorem of Algebra states that every polynomial over the complex numbers factors into distinct linear factors.

Theorem 23.33 Fundamental Theorem of Algebra. *The field of complex numbers is algebraically closed; that is, every polynomial in $\mathbb{C}[x]$ has a root in $\mathbb{C}$.*

PROOF. Suppose that E is a proper finite field extension of the complex numbers. Since any finite extension of a field of characteristic zero is a simple extension, there exists an $\alpha \in E$ such that $E = \mathbb{C}(\alpha)$ with α the root of an irreducible polynomial $f(x)$ in $\mathbb{C}[x]$. The splitting field L of $f(x)$ is a finite normal separable extension of $\mathbb{C}$ that contains E. We must show that it is impossible for L to be a proper extension of $\mathbb{C}$.

Suppose that L is a proper extension of $\mathbb{C}$. Since L is the splitting field of $f(x)(x^2 + 1)$ over $\mathbb{R}$, L is a finite normal separable extension of $\mathbb{R}$. Let K be the fixed field of a Sylow 2-subgroup G of $G(L/\mathbb{R})$. Then $L \supset K \supset \mathbb{R}$ and $|G(L/K)| = [L : K]$. Since $[L : \mathbb{R}] = [L : K][K : \mathbb{R}]$, we know that $[K : \mathbb{R}]$ must be odd. Consequently, $K = \mathbb{R}(\beta)$ with β having a minimal polynomial $f(x)$ of odd degree. Therefore, $K = \mathbb{R}$.

We now know that $G(L/\mathbb{R})$ must be a 2-group. It follows that $G(L/\mathbb{C})$ is a 2-group. We have assumed that $L \neq \mathbb{C}$; therefore, $|G(L/\mathbb{C})| \geq 2$. By the first Sylow Theorem and the Fundamental Theorem of Galois Theory, there exists a subgroup G of $G(L/\mathbb{C})$ of index 2 and a field E fixed elementwise by G. Then $[E : \mathbb{C}] = 2$ and there exists an element $\gamma \in E$ with minimal polynomial $x^2 + bx + c$ in $\mathbb{C}[x]$. This polynomial has roots $(-b \pm \sqrt{b^2 - 4c})/2$ that are in $\mathbb{C}$, since $b^2 - 4c$ is in $\mathbb{C}$. This is impossible; hence, $L = \mathbb{C}$. □

Although our proof was strictly algebraic, we were forced to rely on results from calculus. It is necessary to assume the completeness axiom from analysis to show that every polynomial of odd degree has a real root and that every positive real number has a square root. It seems that there is no possible way to avoid this difficulty and formulate a purely algebraic argument. It is somewhat amazing that there are several elegant proofs of the Fundamental Theorem of Algebra that use complex analysis. It is also interesting to note that we can obtain a proof of such an important theorem from two very different fields of mathematics.

23.4 Exercises

1. Compute each of the following Galois groups. Which of these field extensions are normal field extensions? If the extension is not normal, find a normal extension of $\mathbb{Q}$ in which the extension field is contained.

(a) $G(\mathbb{Q}(\sqrt{30})/\mathbb{Q})$

(b) $G(\mathbb{Q}(\sqrt[4]{5})/\mathbb{Q})$

(c) $G(\mathbb{Q}(\sqrt{2}, \sqrt{3}, \sqrt{5})/\mathbb{Q})$

(d) $G(\mathbb{Q}(\sqrt{2}, \sqrt[3]{2}, i)/\mathbb{Q})$

(e) $G(\mathbb{Q}(\sqrt{6}, i)/\mathbb{Q})$

2. Determine the separability of each of the following polynomials.

(a) $x^3 + 2x^2 - x - 2$ over $\mathbb{Q}$ (c) $x^4 + x^2 + 1$ over $\mathbb{Z}_3$

(b) $x^4 + 2x^2 + 1$ over $\mathbb{Q}$ (d) $x^3 + x^2 + 1$ over $\mathbb{Z}_2$

3.   Give the order and describe a generator of the Galois group of GF(729) over GF(9).

4.   Determine the Galois groups of each of the following polynomials in $\mathbb{Q}[x]$; hence, determine the solvability by radicals of each of the polynomials.

(a) $x^5 - 12x^2 + 2$ (f) $(x^2 - 2)(x^2 + 2)$

(b) $x^5 - 4x^4 + 2x + 2$ (g) $x^8 - 1$

(c) $x^3 - 5$

(d) $x^4 - x^2 - 6$ (h) $x^8 + 1$

(e) $x^5 + 1$ (i) $x^4 - 3x^2 - 10$

5.   Find a primitive element in the splitting field of each of the following polynomials in $\mathbb{Q}[x]$.

(a) $x^4 - 1$ (c) $x^4 - 2x^2 - 15$

(b) $x^4 - 8x^2 + 15$ (d) $x^3 - 2$

6.   Prove that the Galois group of an irreducible quadratic polynomial is isomorphic to $\mathbb{Z}_2$.

7.   Prove that the Galois group of an irreducible cubic polynomial is isomorphic to S_3 or $\mathbb{Z}_3$.

8.   Let $F \subset K \subset E$ be fields. If E is a normal extension of F, show that E must also be a normal extension of K.

9.   Let G be the Galois group of a polynomial of degree n. Prove that $|G|$ divides $n!$.

10.   Let $F \subset E$. If $f(x)$ is solvable over F, show that $f(x)$ is also solvable over E.

11.   Construct a polynomial $f(x)$ in $\mathbb{Q}[x]$ of degree 7 that is not solvable by radicals.

12.   Let p be prime. Prove that there exists a polynomial $f(x) \in \mathbb{Q}[x]$ of degree p with Galois group isomorphic to S_p. Conclude that for each prime p with $p \geq 5$ there exists a polynomial of degree p that is not solvable by radicals.

13.   Let p be a prime and $\mathbb{Z}_p(t)$ be the field of rational functions over $\mathbb{Z}_p$. Prove that $f(x) = x^p - t$ is an irreducible polynomial in $\mathbb{Z}_p(t)[x]$. Show that $f(x)$ is not separable.

14.   Let E be an extension field of F. Suppose that K and L are two intermediate fields. If there exists an element $\sigma \in G(E/F)$ such that $\sigma(K) = L$, then K and L are said to be **conjugate fields.** Prove that K and L are conjugate if and only if $G(E/K)$ and $G(E/L)$ are conjugate subgroups of $G(E/F)$.

15.   Let $\sigma \in \text{Aut}(\mathbb{R})$. If a is a positive real number, show that $\sigma(a) > 0$.

16.   Let K be the splitting field of $x^3 + x^2 + 1 \in \mathbb{Z}_2[x]$. Prove or disprove that K is an extension by radicals.

17.   Let F be a field such that $\text{char}(F) \neq 2$. Prove that the splitting field of $f(x) = ax^2 + bx + c$ is $F(\sqrt{\alpha})$, where $\alpha = b^2 - 4ac$.

18.   Prove or disprove: Two different subgroups of a Galois group will have different fixed fields.

19. Let K be the splitting field of a polynomial over F. If E is a field extension of F contained in K and $[E : F] = 2$, then E is the splitting field of some polynomial in $F[x]$.

20. We know that the cyclotomic polynomial

$$\Phi_p(x) = \frac{x^p - 1}{x - 1} = x^{p-1} + x^{p-2} + \cdots + x + 1$$

is irreducible over $\mathbb{Q}$ for every prime p. Let ω be a zero of $\Phi_p(x)$, and consider the field $\mathbb{Q}(\omega)$.

(a) Show that $\omega, \omega^2, \ldots, \omega^{p-1}$ are distinct zeros of $\Phi_p(x)$, and conclude that they are all the zeros of $\Phi_p(x)$.

(b) Show that $G(\mathbb{Q}(\omega)/\mathbb{Q})$ is abelian of order $p - 1$.

(c) Show that the fixed field of $G(\mathbb{Q}(\omega)/\mathbb{Q})$ is $\mathbb{Q}$.

21. Let F be a finite field or a field of characteristic zero. Let E be a finite normal extension of F with Galois group $G(E/F)$. Prove that $F \subset K \subset L \subset E$ if and only if $\{\text{id}\} \subset G(E/L) \subset G(E/K) \subset G(E/F)$.

22. Let F be a field of characteristic zero and let $f(x) \in F[x]$ be a separable polynomial of degree n. If E is the splitting field of $f(x)$, let $\alpha_1, \ldots, \alpha_n$ be the roots of $f(x)$ in E. Let $\Delta = \prod_{i<j}(\alpha_i - \alpha_j)$. We define the **discriminant** of $f(x)$ to be Δ^2.

(a) If $f(x) = x^2 + bx + c$, show that $\Delta^2 = b^2 - 4c$.

(b) If $f(x) = x^3 + px + q$, show that $\Delta^2 = -4p^3 - 27q^2$.

(c) Prove that Δ^2 is in F.

(d) If $\sigma \in G(E/F)$ is a transposition of two roots of $f(x)$, show that $\sigma(\Delta) = -\Delta$.

(e) If $\sigma \in G(E/F)$ is an even permutation of the roots of $f(x)$, show that $\sigma(\Delta) = \Delta$.

(f) Prove that $G(E/F)$ is isomorphic to a subgroup of A_n if and only if $\Delta \in F$.

(g) Determine the Galois groups of $x^3 + 2x - 4$ and $x^3 + x - 3$.

References and Suggested Readings

[1] Artin, E. *Theory: Lectures Delivered at the University of Notre Dame (Notre Dame Mathematical Lectures, Number 2).* Dover, Mineola, NY, 1997.

[2] Edwards, H. M. *Galois Theory.* Springer-Verlag, New York, 1984.

[3] Fraleigh, J. B. *A First Course in Abstract Algebra.* 7th ed. Pearson, Upper Saddle River, NJ, 2003.

[4] Gaal, L. *Classical Galois Theory with Examples.* American Mathematical Society, Providence, 1979.

[5] Garling, D. J. H. *A Course in Galois Theory.* Cambridge University Press, Cambridge, 1986.

[6] Kaplansky, I. *Fields and Rings.* 2nd ed. University of Chicago Press, Chicago, 1972.

[7] Rothman, T. "The Short Life of Évariste Galois," *Scientific American*, April 1982, 136–49.

23.5   Sage

Again, our competence at examining fields with Sage will allow us to study the main concepts of Galois Theory easily. We will thoroughly examine Example 7 carefully using our computational tools.

Galois Groups

We will repeat Example 23.24 and analyze carefully the splitting field of the polynomial $p(x) = x^4 - 2$. We begin with an initial field extension containing at least one root.

```
x = polygen(QQ, 'x')
N.<a> = NumberField(x^4 - 2); N
```

```
Number Field in a with defining polynomial x^4 - 2
```

The `.galois_closure()` method will create an extension containing all of the roots of the defining polynomial of a number field.

```
L.<b> = N.galois_closure(); L
```

```
Number Field in b with defining polynomial x^8 + 28*x^4 + 2500
```

```
L.degree()
```

```
8
```

```
y = polygen(L, 'y')
(y^4 - 2).factor()
```

```
(y - 1/120*b^5 -  19/60*b) *
(y - 1/240*b^5 + 41/120*b) *
(y + 1/240*b^5 - 41/120*b) *
(y + 1/120*b^5 +  19/60*b)
```

From the factorization, it is clear that `L` is the splitting field of the polynomial, even if the factorization is not pretty. It is easy to then obtain the Galois group of this field extension.

```
G = L.galois_group(); G
```

```
Galois group of Number Field in b with
defining polynomial x^8 + 28*x^4 + 2500
```

We can examine this group, and identify it. Notice that since the field is a degree 8 extension, the group is described as a permutation group on 8 symbols. (It is just a coincidence that the group has 8 elements.) With a paucity of nonabelian groups of order 8, it is not hard to guess the nature of the group.

```
G.is_abelian()
```

```
False
```

```
G.order()
```

```
8
```

```
G.list()
```

```
[(), (1,2,8,7)(3,4,6,5),
(1,3)(2,5)(4,7)(6,8), (1,4)(2,3)(5,8)(6,7),
(1,5)(2,6)(3,7)(4,8), (1,6)(2,4)(3,8)(5,7),
(1,7,8,2)(3,5,6,4), (1,8)(2,7)(3,6)(4,5)]
```

```
G.is_isomorphic(DihedralGroup(4))
```

```
True
```

That's it. But maybe not very satisfying. Let us dig deeper for more understanding. We will start over and create the splitting field of $p(x) = x^4 - 2$ again, but the primary difference is that we will make the roots extremely obvious so we can work more carefully with the Galois group and the fixed fields. Along the way, we will see another example of linear algebra enabling certain computations. The following construction should be familiar by now.

```
x = polygen(QQ, 'x')
p = x^4 - 2
N.<a> = NumberField(p); N
```

```
Number Field in a with defining polynomial x^4 - 2
```

```
y = polygen(N, 'y')
p = p.subs(x=y)
p.factor()
```

```
(y - a) * (y + a) * (y^2 + a^2)
```

```
M.<b> = NumberField(y^2 + a^2); M
```

```
Number Field in b with defining polynomial y^2 + a^2 over
its base field
```

```
z = polygen(M, 'z')
(z^4 - 2).factor()
```

```
(z - b) * (z - a) * (z + a) * (z + b)
```

The important thing to notice here is that we have arranged the splitting field so that the four roots, `a`, `-a`, `b`, `-b`, are very simple functions of the generators. In more traditional notation, `a` is $2^{\frac{1}{4}} = \sqrt[4]{2}$, and `b` is $2^{\frac{1}{4}}i = \sqrt[4]{2}i$ (or their negatives).

We will find it easier to compute in the flattened tower, a now familiar construction.

```
L.<c> = M.absolute_field(); L
```

```
Number Field in c with defining polynomial x^8 + 28*x^4 + 2500
```

```
fromL, toL = L.structure()
```

We can return to our original polynomial (over the rationals), and ask for its roots in the flattened tower, custom-designed to contain these roots.

```
roots = p.roots(ring=L, multiplicities=False); roots
```

```
[1/120*c^5 +  19/60*c,
 1/240*c^5 - 41/120*c,
-1/240*c^5 + 41/120*c,
-1/120*c^5 -  19/60*c]
```

Hmmm. Do those look right? If you look back at the factorization obtained in the field constructed with the `.galois_closure()` method, then they look right. But we can do better.

```
[fromL(r) for r in roots]
```

```
[b, a, -a, -b]
```

Yes, those are the roots.

The `End()` command will create the group of automorphisms of the field `L`.

```
G = End(L); G
```

```
Automorphism group of Number Field in c with
defining polynomial x^8 + 28*x^4 + 2500
```

We can check that each of these automorphisms fixes the rational numbers elementwise. If a field homomorphism fixes 1, then it will fix the integers, and thus fix all fractions of integers.

```
[tau(1) for tau in G]
```

```
[1, 1, 1, 1, 1, 1, 1, 1]
```

So each element of `G` fixes the rationals elementwise and thus `G` is the Galois group of the splitting field `L` over the rationals.

Proposition 23.5 is fundamental. It says every automorphism in the Galois group of a field extension creates a permutation of the roots of a polynomial with coefficients in the base field. We have all of those ingredients here. So we will evaluate each automorphism of the Galois group at each of the four roots of our polynomial, which in each case should be another root. (We use the `Sequence()` constructor just to get nicely-aligned output.)

```
Sequence([[fromL(tau(r)) for r in roots] for tau in G], cr=True)
```

```
[
[b, a, -a, -b],
[-b, -a, a, b],
[a, -b, b, -a],
[b, -a, a, -b],
[-a, -b, b, a],
[a, b, -b, -a],
[-b, a, -a, b],
[-a, b, -b, a]
]
```

Each row of the output is a list of the roots, but permuted, and so corresponds to a permutation of four objects (the roots). For example, the second row shows the second automorphism interchanging `a` with `-a`, and `b` with `-b`. (Notice that the first row is the result of the identity automorphism, so we can mentally comine the first row with any other row to imagine a "two-row" form of a permutation.) We can number the roots, 1 through 4, and create each permutation as an element of S_4. It is overkill, but we can then build the permutation group by letting *all* of these elements generate a group.

```
S4 = SymmetricGroup(4)
elements = [S4([1, 2, 3, 4]),
            S4([4, 3, 2, 1]),
            S4([2, 4, 1, 3]),
            S4([1, 3, 2, 4]),
```

```
              S4([3, 4, 1, 2]),
              S4([2, 1, 4, 3]),
              S4([4, 2, 3, 1]),
              S4([3, 1, 4, 2])]
elements
```

```
[(), (1,4)(2,3), (1,2,4,3), (2,3), (1,3)(2,4),
 (1,2)(3,4), (1,4), (1,3,4,2)]
```

```
P = S4.subgroup(elements)
P.is_isomorphic(DihedralGroup(4))
```

```
True
```

Notice that we now have built an isomorphism from the Galois group to a group of permutations *using just four symbols*, rather than the eight used previously.

Fixed Fields

In a previous Sage exercise, we computed the fixed fields of single field automorphisms for finite fields. This was "easy" in the sense that we could just test every element of the field to see if it was fixed, since the field was finite. Now we have an infinite field extension. How are we going to determine which elements are fixed by individual automorphisms, or subgroups of automorphisms?

The answer is to use the vector space structure of the flattened tower. As a degree 8 extension of the rationals, the first 8 powers of the primitive element `c` form a basis when the field is viewed as a vector space with the rationals as the scalars. It is sufficient to know how each field automorphism behaves on this basis to fully specify the definition of the automorphism. To wit,

$$
\begin{aligned}
\tau(x) = \tau\left(\sum_{i=0}^{7} q_i c^i\right) && q_i \in \mathbb{Q} \\
= \sum_{i=0}^{7} \tau(q_i)\tau(c^i) && \tau \text{ is a field automorphism} \\
= \sum_{i=0}^{7} q_i \tau(c^i) && \text{rationals are fixed}
\end{aligned}
$$

So we can compute the value of a field automorphism at any linear combination of powers of the primitive element as a linear combination of the values of the field automorphism at just the powers of the primitive element. This is known as the "power basis", which we can obtain simply with the `.power_basis()` method. We will begin with an example of how we can use this basis. We will illustrate with the fourth automorphism of the Galois group. Notice that the `.vector()` method is a convenience that strips a linear combination of the powers of `c` into a vector of just the coefficients. (Notice too that τ is totally defined by the value of $\tau(c)$, since as a field automorphism $\tau(c^k) = (\tau(c))^k$. However, we still need to work with the entire power basis to exploit the vector space structure.)

```
basis = L.power_basis(); basis
```

```
[1, c, c^2, c^3, c^4, c^5, c^6, c^7]
```

```
tau = G[3]
z = 4 + 5*c+ 6*c^3-7*c^6
tz = tau(4 + 5*c+ 6*c^3-7*c^6); tz
```

```
11/250*c^7 - 98/25*c^6 + 1/12*c^5 + 779/125*c^3 +
6006/25*c^2 - 11/6*c + 4
```

```
tz.vector()
```

```
(4, -11/6, 6006/25, 779/125, 0, 1/12, -98/25, 11/250)
```

```
tau_matrix = column_matrix([tau(be).vector() for be in basis])
tau_matrix
```

```
[  1        0        0        0    -28        0        0          0]
[  0   -11/30        0        0      0   779/15        0          0]
[  0        0   -14/25        0      0        0   -858/25         0]
[  0        0        0   779/750     0        0        0   -4031/375]
[  0        0        0        0     -1        0        0          0]
[  0     1/60        0        0      0    11/30        0          0]
[  0        0    -1/50        0      0        0    14/25          0]
[  0        0        0   11/1500     0        0        0    -779/750]
```

```
tau_matrix*z.vector()
```

```
(4, -11/6, 6006/25, 779/125, 0, 1/12, -98/25, 11/250)
```

```
tau_matrix*(z.vector()) == (tau(z)).vector()
```

```
True
```

The last line expresses the fact that `tau_matrix` is a matrix representation of the field automorphism, viewed as a linear transformation of the vector space structure. As a representation of an invertible field homomorphism, the matrix is invertible. As an order 2 permutation of the roots, the inverse of the matrix is itself. But these facts are just verifications that we have the right thing, we are interested in other properties.

To construct fixed fields, we want to find elements fixed by automorphisms. Continuing with `tau` from above, we seek elements `z` (written as vectors) such that `tau_matrix*z=z`. These are eigenvectors for the eigenvalue 1, or elements of the null space of (`tau_matrix - I`) (null spaces are obtained with `.right_kernel()` in Sage).

```
K = (tau_matrix-identity_matrix(8)).right_kernel(); K
```

```
Vector space of degree 8 and dimension 4 over Rational Field
Basis matrix:
[    1     0     0     0     0     0     0      0]
[    0     1     0     0     0   1/38    0      0]
[    0     0     1     0     0     0  -1/22     0]
[    0     0     0     1     0     0     0  1/278]
```

Each row of the basis matrix is a vector representing an element of the field, specifically 1, `c + (1/38)*c^5`, `c^2 - (1/22)*c^6`, `c^3 + (1/278)*c^7`. Let's take a closer look at these fixed elements, in terms we recognize.

```
fromL(1)
```

1

```
fromL(c + (1/38)*c^5)
```

60/19*b

```
fromL(c^2 - (1/22)*c^6)
```

150/11*a^2

```
fromL(c^3 + (1/278)*c^7)
```

1500/139*a^2*b

Any element fixed by `tau` will be a linear combination of these four elements. We can ignore any rational multiples present, the first element is just saying the rationals are fixed, and the last element is just a product of the middle two. So fundamentally `tau` is fixing rationals, `b` (which is $\sqrt[4]{2}i$) and `a^2` (which is $\sqrt{2}$). Furthermore, `b^2` = `-a^2` (the check follows), so we can create any fixed element of `tau` by just adjoining `b=`$\sqrt[4]{2}i$ to the rationals. So the elements fixed by `tau` are $\mathbb{Q}(\sqrt[4]{2}i)$.

```
a^2 + b^2
```

0

Galois Correspondence

The entire subfield structure of our splitting field is determined by the subgroup structure of the Galois group (Theorem 23.22), which is isomorphic to a group we know well. What are the subgroups of our Galois group, expressed as permutation groups? (For brevity, we just list the *generators* of each subgroup.)

```
sg = P.subgroups();
[H.gens() for H in sg]
```

```
[[()],
 [(2,3)],
 [(1,4)],
 [(1,4)(2,3)],
 [(1,2)(3,4)],
 [(1,3)(2,4)],
 [(2,3), (1,4)],
 [(1,2)(3,4), (1,4)(2,3)],
 [(1,3,4,2), (1,4)(2,3)],
 [(2,3), (1,2)(3,4), (1,4)]]
```

```
[H.order() for H in sg]
```

[1, 2, 2, 2, 2, 2, 4, 4, 4, 8]

`tau` above is the fourth element of the automorphism group, and the fourth permutation in `elements` is the permutation (2,3), the generator (of order 2) for the second subgroup. So as the only nontrivial element of this subgroup, we know that the corresponding fixed field is $\mathbb{Q}(\sqrt[4]{2}i)$.

Let us analyze another subgroup of order 2, without all the explanation, and starting with the subgroup. The sixth subgroup is generated by the fifth automorphism, so let us determine the elements that are fixed.

```
tau = G[4]
tau_matrix = column_matrix([tau(be).vector() for be in basis])
(tau_matrix-identity_matrix(8)).right_kernel()
```

```
Vector space of degree 8 and dimension 4 over Rational Field
Basis matrix:
[      1       0       0       0       0       0       0       0]
[      0       1       0       0       0   1/158       0       0]
[      0       0       1       0       0       0    1/78       0]
[      0       0       0       1       0       0       0  13/614]
```

```
fromL(tau(1))
```

1

```
fromL(tau(c+(1/158)*c^5))
```

120/79*b - 120/79*a

```
fromL(tau(c^2+(1/78)*c^6))
```

-200/39*a*b

```
fromL(tau(c^3+(13/614)*c^7))
```

3000/307*a^2*b + 3000/307*a^3

The first element indicates that the rationals are fixed (we knew that). Scaling the second element gives b - a as a fixed element. Scaling the third and fourth fixed elements, we recognize that they can be obtained from powers of b - a.

```
(b-a)^2
```

-2*a*b

```
(b-a)^3
```

2*a^2*b + 2*a^3

So the fixed field of this subgroup can be formed by adjoining b - a to the rationals, which in mathematical notation is $\sqrt[4]{2}i - \sqrt[4]{2} = (1-i)\sqrt[4]{2}$, so the fixed field is $\mathbb{Q}(\sqrt[4]{2}i - \sqrt[4]{2}) = \mathbb{Q}((1-i)\sqrt[4]{2})$.

We can create this fixed field, though as created here it is not strictly a subfield of L. We will use an expression for b - a that is a linear combination of powers of c.

```
subinfo = L.subfield((79/120)*(c+(1/158)*c^5)); subinfo
```

```
(Number Field in c0 with defining polynomial x^4 + 8, Ring morphism:
  From: Number Field in c0 with defining polynomial x^4 + 8
  To:   Number Field in c with defining polynomial x^8 + 28*x^4 + 2500
  Defn: c0 |--> 1/240*c^5 + 79/120*c)
```

The .subfield() method returns a pair. The first item is a new number field, isomorphic to a subfield of L. The second item is an injective mapping from the new number field into L. In this case, the image of the primitive element c0 is the element we have specified as the

generator of the subfield. The primitive element of the new field will satisfy the defining polynomial $x^4 + 8$ — you can check that $(1-i)\sqrt[4]{2}$ is indeed a root of the polynomial $x^4 + 8$.

There are five subgroups of order 2, we have found fixed fields for two of them. The other three are similar, so it would be a good exercise to work through them. Our automorphism group has three subgroups of order 4, and at least one of each possible type (cyclic versus non-cyclic). Fixed fields of larger subgroups require that we find elements fixed by all of the automorphisms in the subgroup. (We were conveniently ignoring the identity automorphism above.) This will require more computation, but will restrict the possibilities (smaller fields) to where it will be easier to deduce a primitive element for each field.

The seventh subgroup is generated by two elements of order 2 and is composed entirely of elements of order 2 (except the identity), so is isomorphic to $\mathbb{Z}_2 \times \mathbb{Z}_2$. The permutations correspond to automorphisms number 0, 1, 3, and 6. To determine the elements fixed by *all four* automorphisms, we will build the kernel for each one and as we go, we form the *intersection* of all four kernels. We will work via a loop over the four automorphisms.

```
V = QQ^8
for tau in [G[0], G[1], G[3], G[6]]:
  tau_matrix = column_matrix([tau(be).vector() for be in basis])
  K = (tau_matrix-identity_matrix(8)).right_kernel()
  V = V.intersection(K)
V
```

```
Vector space of degree 8 and dimension 2 over Rational Field
Basis matrix:
[    1     0     0     0     0     0     0     0]
[    0     0     1     0     0     0 -1/22     0]
```

Outside of the rationals, there is a single fixed element.

```
fromL(tau(c^2 - (1/22)*c^6))
```

```
150/11*a^2
```

Removing a scalar multiple, our primitive element is `a^2`, which mathematically is $\sqrt{2}$, so the fixed field is $\mathbb{Q}(\sqrt{2})$. Again, we can build this fixed field, but ignore the mapping.

```
F, mapping = L.subfield((11/150)*(c^2 - (1/22)*c^6))
F
```

```
Number Field in c0 with defining polynomial x^2 - 2
```

One more subgroup. The penultimate subgroup has a permutation of order 4 as a generator, so is a cyclic group of order 4. The individual permutations of the subgroup correspond to automorphisms 0, 1, 2, 7.

```
V = QQ^8
for tau in [G[0], G[1], G[2], G[7]]:
  tau_matrix = column_matrix([tau(be).vector() for be in basis])
  K = (tau_matrix-identity_matrix(8)).right_kernel()
  V = V.intersection(K)
V
```

```
Vector space of degree 8 and dimension 2 over Rational Field
Basis matrix:
[1 0 0 0 0 0 0 0]
[0 0 0 0 1 0 0 0]
```

So we compute the primitive element.

```
fromL(tau(c^4))
```

```
-24*a^3*b - 14
```

Since rationals are fixed, we can remove the -14 and the multiple and take `a^3*b` as the primitive element. Mathematically, this is $2i$, so we might as well use just i as the primitive element and the fixed field is $\mathbb{Q}(i)$. We can then build the fixed field (and ignore the mapping also returned).

```
F, mapping = L.subfield((c^4+14)/-48)
F
```

```
Number Field in c0 with defining polynomial x^2 + 1
```

There is one more subgroup of order 4, which we will leave as an exercise to analyze. There are also two trivial subgroups (the identity and the full group) which are not very interesting or surprising.

If the above seems like too much work, you can always just have Sage do it all with the `.subfields()` method.

```
L.subfields()
```

```
[
(Number Field in c0 with defining polynomial x,
 Ring morphism:
   From: Number Field in c0 with defining polynomial x
   To:   Number Field in c with defining polynomial x^8 + 28*x^4 + 2500
   Defn: 0 |--> 0,
 None),
(Number Field in c1 with defining polynomial x^2 + 112*x + 40000,
 Ring morphism:
   From: Number Field in c1 with defining polynomial x^2 + 112*x + 40000
   To:   Number Field in c with defining polynomial x^8 + 28*x^4 + 2500
   Defn: c1 |--> 4*c^4,
 None),
(Number Field in c2 with defining polynomial x^2 + 512,
 Ring morphism:
   From: Number Field in c2 with defining polynomial x^2 + 512
   To:   Number Field in c with defining polynomial x^8 + 28*x^4 + 2500
   Defn: c2 |--> 1/25*c^6 + 78/25*c^2,
 None),
(Number Field in c3 with defining polynomial x^2 - 288,
 Ring morphism:
   From: Number Field in c3 with defining polynomial x^2 - 288
   To:   Number Field in c with defining polynomial x^8 + 28*x^4 + 2500
   Defn: c3 |--> -1/25*c^6 + 22/25*c^2,
 None),
(Number Field in c4 with defining polynomial x^4 + 112*x^2 + 40000,
 Ring morphism:
   From: Number Field in c4 with defining polynomial x^4 + 112*x^2 +
       40000
   To:   Number Field in c with defining polynomial x^8 + 28*x^4 + 2500
   Defn: c4 |--> 2*c^2,
 None),
(Number Field in c5 with defining polynomial x^4 + 648,
 Ring morphism:
   From: Number Field in c5 with defining polynomial x^4 + 648
```

```
   To:   Number Field in c with defining polynomial x^8 + 28*x^4 + 2500
   Defn: c5 |--> 1/80*c^5 + 79/40*c,
 None),
(Number Field in c6 with defining polynomial x^4 + 8,
Ring morphism:
  From: Number Field in c6 with defining polynomial x^4 + 8
  To:   Number Field in c with defining polynomial x^8 + 28*x^4 + 2500
  Defn: c6 |--> -1/80*c^5 + 1/40*c,
  None),
(Number Field in c7 with defining polynomial x^4 - 512,
 Ring morphism:
   From: Number Field in c7 with defining polynomial x^4 - 512
   To:   Number Field in c with defining polynomial x^8 + 28*x^4 + 2500
   Defn: c7 |--> -1/60*c^5 + 41/30*c,
 None),
(Number Field in c8 with defining polynomial x^4 - 32,
 Ring morphism:
   From: Number Field in c8 with defining polynomial x^4 - 32
   To:   Number Field in c with defining polynomial x^8 + 28*x^4 + 2500
   Defn: c8 |--> 1/60*c^5 + 19/30*c,
 None),
(Number Field in c9 with defining polynomial x^8 + 28*x^4 + 2500,
 Ring morphism:
   From: Number Field in c9 with defining polynomial x^8 + 28*x^4 + 2500
   To:   Number Field in c with defining polynomial x^8 + 28*x^4 + 2500
   Defn: c9 |--> c,
 Ring morphism:
   From: Number Field in c with defining polynomial x^8 + 28*x^4 + 2500
   To:   Number Field in c9 with defining polynomial x^8 + 28*x^4 + 2500
   Defn: c |--> c9)
]
```

Ten subfields are described, which is what we would expect, given the 10 subgroups of the Galois group. Each begins with a new number field that is a subfield. Technically, each is not a subset of L, but the second item returned for each subfield is an injective homomorphism, also known generally as an "embedding." Each embedding describes how a primitive element of the subfield translates to an element of L. Some of these primitive elements could be manipulated (as we have done above) to yield slightly simpler minimal polynomials, but the results are quite impressive nonetheless. Each item in the list has a third component, which is almost always None, except when the subfield is the whole field, and then the third component is an injective homomorphism "in the other direction."

Normal Extensions

Consider the third subgroup in the list above, generated by the permutation (1,4). As a subgroup of order 2, it only has one nontrivial element, which here corresponds to the seventh automorphism. We determine the fixed elements as before.

```
tau = G[6]
tau_matrix = column_matrix([tau(be).vector() for be in basis])
(tau_matrix-identity_matrix(8)).right_kernel()
```

```
Vector space of degree 8 and dimension 4 over Rational Field
Basis matrix:
[    1    0    0    0    0     0    0    0]
[    0    1    0    0    0 -1/82    0    0]
```

```
[     0      0      1      0      0      0  -1/22       0]
[     0      0      0      1      0      0       0  11/58]
```

```
fromL(tau(1))
```

```
1
```

```
fromL(tau(c+(-1/82)*c^5))
```

```
-120/41*a
```

```
fromL(tau(c^2+(-1/22)*c^6))
```

```
150/11*a^2
```

```
fromL(tau(c^3+(11/58)*c^7))
```

```
3000/29*a^3
```

As usual, ignoring rational multiples, we see powers of a and recognize that a alone will be a primitive element for the fixed field, which is thus $\mathbb{Q}(\sqrt[4]{2})$. Recognize that a was our first root of $x^4 - 2$, and was used to create the first part of original tower, N. So N is both $\mathbb{Q}(\sqrt[4]{2})$ and the fixed field of $H = \langle (1, 4) \rangle$.

$\mathbb{Q}(\sqrt[4]{2})$ contains at least one root of the irreducible $x^4 - 2$, but not all of the roots (witness the factorization above) and therefore does not qualify as a normal extension. By part (4) of Theorem 23.22 the automorphism group of the extension is not normal in the full Galois group.

```
sg[2].is_normal(P)
```

```
False
```

As expected.

23.6   Sage Exercises

1.   In the analysis of Example 23.24 with Sage, two subgroups of order 2 and one subgroup of order 4 were not analyzed. Determine the fixed fields of these three subgroups.

2.   Build the splitting field of $p(x) = x^3 - 6x^2 + 12x - 10$ and then determine the Galois group of $p(x)$ as a concrete group of explicit permutations. Build the lattice of subgroups of the Galois group, again using the same explicit permutations. Using the Fundamental Theorem of Galois Theory, construct the subfields of the splitting field. Include your supporting documentation in your submitted Sage worksheet. Also, submit a written component of this assignment containing a complete layout of the subgroups and subfields, written entirely with mathematical notation and with no Sage commands, designed to illustrate the correspondence between the two. All you need here is the graphical layout, suitably labeled — the Sage worksheet will substantiate your work.

3.   The polynomial $x^5 - x - 1$ has all of the symmetric group S_5 as its Galois group. Because S_5 is not solvable, we know this polynomial to be an example of a quintic polynomial that is not solvable by radicals. Unfortunately, asking Sage to compute this Galois group takes far too long. So this exercise will simulate that experience with a slightly smaller example.
     Consider the polynomial $p(x) = x^4 + x + 1$.

(a) Build the splitting field of $p(x)$ one root at a time. Create an extension, factor there, discard linear factors, use the remaining irreducible factor to extend once more. Repeat until $p(x)$ factors completely. Be sure to do a final extension via just a linear factor. This is a little silly, and Sage will seem to ignore your final generator (so you will want to setermine what it is equivalent to in terms of the previous gfenerators). Directions below depend on taking this extra step.

(b) Factor the original polynomial over the final extension field in the tower. What is boring about this factorization in comparison to some other examples we have done?

(c) Construct the full tower as an absolute field over $\mathbb{Q}$. From the degree of this extension and the degree of the original polynomial, infer the Galois group of the polynomial.

(d) Using the mappings that allow you to translate between the tower and the absolute field (obtained from the `.structure()` method), choose one of the roots (any one) and express it in terms of the single generator of the absolute field. Then reverse the procedure and express the single generator of the absolute field in terms of the roots in the tower.

(e) Compute the group of automorphisms of the absolute field (but don't display the whole group in what you submit). Take all four roots (including your silly one from the last step of the tower construction) and apply each field automorphism to the four roots (creating the guaranteed permutations of the roots). Comment on what you see.

(f) There is one nontrivial automorphism that has an especially simple form (it is the second one for me) when applied to the generator of the absolute field. What does this automorphism do to the roots of $p(x)$?

(g) Consider the extension of $\mathbb{Q}$ formed by adjoining just one of the roots. This is a subfield of the splitting field of the polynomial, so is the fixed field of a subgroup of the Galois group. Give a simple description of the corresponding subgroup using language we typically only apply to permutation groups.

4. Return to the splitting field of the quintic discussed in the introduction to the previous problem ($x^5 - x - 1$). Create the first two intermediate fields by adjoining two roots (one at a time). But instead of factoring at each step to get a new irreducible polynomial, *divide* by the linear factor you *know* is a factor. In general, the quotient might factor further, but in this exercise presume it does not. In other words, act as if your quotient by the linear factor is irreducible. If it is not, then the `NumberField()` command should complain (which it will not).

After adjoining two roots, create the extension producing a third root, and do the division. You should now have a quadratic factor. Assuming the quadratic is irreducible (it is) argue that you have enough evidence to establish the order of the Galois group, and hence can determine *exactly* which group it is.

You can try to use this quadratic factor to create one more step in the extensions, and you will arrive at the splitting field, as can be seen with logic or division. However, this could take a long time to complete (save your work beforehand!). You can try passing the `check=False` argument to the `NumberField()` command — this will bypass checking irreducibility.

5. Create the finite field of order 3^6, letting Sage supply the default polynomial for its construction. The polynomial $x^6 + x^2 + 2x + 1$ is irreducible over this finite field. Check

that this polynomial splits in the finite field, and then use the `.roots()` method to collect the roots of the polynomial. Get the group of automorphisms of the field with the `End()` command.

You now have all of the pieces to associate each field automorphism with a permutation of the roots. From this, identify the Galois group and all of its subgroups. For each subgroup, determine the fixed field. You might find the roots easier to work with if you use the `.log()` method to identify them as powers of the field's multiplicative generator.

Your Galois group in this example will be abelian. So every subgroup is normal, and hence any extension is also normal. Can you extend this example by choosing a nontrivial intermediate field with a nontrivial irreducible polynomial that has all of its roots in the intermediate field and a nontrivial irreducible polynomial with none of its roots in the intermediate field?

Your results here are "typical" in the sense that the particular field or irreducible polynomial makes little difference in the qualitative nature of the results.

6.   The splitting field for the irreducible polynomial $p(x) = x^7 - 7x + 3$ has degree 168 (hence this is the order of the Galois group). This polynomial is derived from an "Elkies trinomial curve," a hyperelliptic curve (below) that produces polynomials with interesting Galois groups:

$$y^2 = x(81x^5 + 396x^4 + 738x^3 + 660x^2 + 269x + 48)$$

For $p(x)$ the resulting Galois group is $PSL(2,7)$, a simple group. If $SL(2,7)$ is all 2×2 matrices over $\mathbb{Z}_7$ with determinant 1, then $PSL(2,7)$ is the quotient by the subgroup $\{I_2, -I_2\}$. It is the second-smallest non-abelian simple group (after A_5).

See how far you can get in using Sage to build this splitting field. A degree 7 extension will yield one linear factor, and a subsequent degree 6 extension will yield two linear factors, leaving a quartic factor. Here is where the computations begin to slow down. If we believe that the splitting field has degree 168, then we know that adding a root from this degree 4 factor will get us to the splitting field. Creating this extension may be possible computationally, but verifying that the quartic splits into linear factors here seems to be infeasible.

7.   Return to Example 23.24, and the complete list of subfields obtainable from the `.subfields()` method applied to the flattened tower. As mentioned, these are technically not subfields, but do have embeddings into the tower. Given two subfields, their respective primitive elements are embedded into the tower, with an image that is a linear combination of powers of the primitive element for the tower.

If one subfield is contained in the other, then the image of the primitive element for the smaller field should be a linear combination of the (appropriate) powers of the image of the primitive element for the larger field. This is a linear algebra computation that should be possible in the tower, relative to the power basis for the whole tower.

Write a procedure to determine if two subfields are related by one being a subset of the other. Then use this procedure to create the lattice of subfields. The eventual goal would be a graphical display of the lattice, using the existing plotting facilities available for lattices, similar to the top half of Figure 23.25. This is a "challenging" exercise, which is code for "it is speculative and has not been tested."

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Hints and Solutions to Selected Exercises

1 Preliminaries

1.3 Exercises

1.3.1. **Hint**. (a) $A \cap B = \{2\}$; (b) $B \cap C = \{5\}$.

1.3.2. **Hint**. (a) $A \times B = \{(a, 1), (a, 2), (a, 3), (b, 1), (b, 2), (b, 3), (c, 1), (c, 2), (c, 3)\}$; (d) $A \times D = \emptyset$.

1.3.6. **Hint**. If $x \in A \cup (B \cap C)$, then either $x \in A$ or $x \in B \cap C$. Thus, $x \in A \cup B$ and $A \cup C$. Hence, $x \in (A \cup B) \cap (A \cup C)$. Therefore, $A \cup (B \cap C) \subset (A \cup B) \cap (A \cup C)$. Conversely, if $x \in (A \cup B) \cap (A \cup C)$, then $x \in A \cup B$ and $A \cup C$. Thus, $x \in A$ or x is in both B and C. So $x \in A \cup (B \cap C)$ and therefore $(A \cup B) \cap (A \cup C) \subset A \cup (B \cap C)$. Hence, $A \cup (B \cap C) = (A \cup B) \cap (A \cup C)$.

1.3.10. **Hint**. $(A \cap B) \cup (A \setminus B) \cup (B \setminus A) = (A \cap B) \cup (A \cap B') \cup (B \cap A') = [A \cap (B \cup B')] \cup (B \cap A') = A \cup (B \cap A') = (A \cup B) \cap (A \cup A') = A \cup B$.

1.3.14. **Hint**. $A \setminus (B \cup C) = A \cap (B \cup C)' = (A \cap A) \cap (B' \cap C') = (A \cap B') \cap (A \cap C') = (A \setminus B) \cap (A \setminus C)$.

1.3.17. **Hint**. (a) Not a map since $f(2/3)$ is undefined; (b) this is a map; (c) not a map, since $f(1/2) = 3/4$ but $f(2/4) = 3/8$; (d) this is a map.

1.3.18. **Hint**. (a) f is one-to-one but not onto. $f(\mathbb{R}) = \{x \in \mathbb{R} : x > 0\}$. (c) f is neither one-to-one nor onto. $f(\mathbb{R}) = \{x : -1 \leq x \leq 1\}$.

1.3.20. **Hint**. (a) $f(n) = n + 1$.

1.3.22. **Hint**. (a) Let $x, y \in A$. Then $g(f(x)) = (g \circ f)(x) = (g \circ f)(y) = g(f(y))$. Thus, $f(x) = f(y)$ and $x = y$, so $g \circ f$ is one-to-one. (b) Let $c \in C$, then $c = (g \circ f)(x) = g(f(x))$ for some $x \in A$. Since $f(x) \in B$, g is onto.

1.3.23. **Hint**. $f^{-1}(x) = (x + 1)/(x - 1)$.

1.3.24. **Hint**. (a) Let $y \in f(A_1 \cup A_2)$. Then there exists an $x \in A_1 \cup A_2$ such that $f(x) = y$. Hence, $y \in f(A_1)$ or $f(A_2)$. Therefore, $y \in f(A_1) \cup f(A_2)$. Consequently, $f(A_1 \cup A_2) \subset f(A_1) \cup f(A_2)$. Conversely, if $y \in f(A_1) \cup f(A_2)$, then $y \in f(A_1)$ or $f(A_2)$. Hence, there exists an x in A_1 or A_2 such that $f(x) = y$. Thus, there exists an $x \in A_1 \cup A_2$ such that $f(x) = y$. Therefore, $f(A_1) \cup f(A_2) \subset f(A_1 \cup A_2)$, and $f(A_1 \cup A_2) = f(A_1) \cup f(A_2)$.

1.3.25. **Hint**. (a) The relation fails to be symmetric. (b) The relation is not reflexive, since 0 is not equivalent to itself. (c) The relation is not transitive.

1.3.28. **Hint**. Let $X = \mathbb{N} \cup \{\sqrt{2}\}$ and define $x \sim y$ if $x + y \in \mathbb{N}$.

2 The Integers

2.3 Exercises

2.3.1. **Hint.** The base case, $S(1) : [1(1+1)(2(1)+1)]/6 = 1 = 1^2$ is true. Assume that $S(k) : 1^2 + 2^2 + \cdots + k^2 = [k(k+1)(2k+1)]/6$ is true. Then

$$1^2 + 2^2 + \cdots + k^2 + (k+1)^2 = [k(k+1)(2k+1)]/6 + (k+1)^2$$
$$= [(k+1)((k+1)+1)(2(k+1)+1)]/6,$$

and so $S(k+1)$ is true. Thus, $S(n)$ is true for all positive integers n.

2.3.3. **Hint.** The base case, $S(4) : 4! = 24 > 16 = 2^4$ is true. Assume $S(k) : k! > 2^k$ is true. Then $(k+1)! = k!(k+1) > 2^k \cdot 2 = 2^{k+1}$, so $S(k+1)$ is true. Thus, $S(n)$ is true for all positive integers n.

2.3.8. **Hint.** Follow the proof in Example 2.4.

2.3.11. **Hint.** The base case, $S(0) : (1+x)^0 - 1 = 0 \geq 0 = 0 \cdot x$ is true. Assume $S(k) : (1+x)^k - 1 \geq kx$ is true. Then

$$(1+x)^{k+1} - 1 = (1+x)(1+x)^k - 1$$
$$= (1+x)^k + x(1+x)^k - 1$$
$$\geq kx + x(1+x)^k$$
$$\geq kx + x$$
$$= (k+1)x,$$

so $S(k+1)$ is true. Therefore, $S(n)$ is true for all positive integers n.

2.3.17. Fibonacci Numbers. **Hint.** For (a) and (b) use mathematical induction. (c) Show that $f_1 = 1$, $f_2 = 1$, and $f_{n+2} = f_{n+1} + f_n$. (d) Use part (c). (e) Use part (b) and Exercise 2.3.16.

2.3.19. **Hint.** Use the Fundamental Theorem of Arithmetic.

2.3.23. **Hint.** Use the Principle of Well-Ordering and the division algorithm.

2.3.27. **Hint.** Since $\gcd(a, b) = 1$, there exist integers r and s such that $ar + bs = 1$. Thus, $acr + bcs = c$.

2.3.29. **Hint.** Every prime must be of the form 2, 3, $6n + 1$, or $6n + 5$. Suppose there are only finitely many primes of the form $6k + 5$.

3 Groups

3.4 Exercises

3.4.1. **Hint.** (a) $3 + 7\mathbb{Z} = \{\ldots, -4, 3, 10, \ldots\}$; (c) $18 + 26\mathbb{Z}$; (e) $5 + 6\mathbb{Z}$.

3.4.2. **Hint.** (a) Not a group; (c) a group.

3.4.6. **Hint.**

$\cdot$	1	5	7	11
1	1	5	7	11
5	5	1	11	7
7	7	11	1	5
11	11	7	5	1

3.4.8.   **Hint**.   Pick two matrices. Almost any pair will work.

3.4.15.   **Hint**.   There is a nonabelian group containing six elements.

3.4.16.   **Hint**.   Look at the symmetry group of an equilateral triangle or a square.

3.4.17.   **Hint**.   The are five different groups of order 8.

3.4.18.   **Hint**.   Let

$$\sigma = \begin{pmatrix} 1 & 2 & \cdots & n \\ a_1 & a_2 & \cdots & a_n \end{pmatrix}$$

be in S_n. All of the a_is must be distinct. There are n ways to choose a_1, $n-1$ ways to choose a_2, ..., 2 ways to choose a_{n-1}, and only one way to choose a_n. Therefore, we can form σ in $n(n-1)\cdots 2 \cdot 1 = n!$ ways.

3.4.25.   **Hint**.

$$\begin{aligned}
(aba^{-1})^n &= (aba^{-1})(aba^{-1})\cdots(aba^{-1}) \\
&= ab(aa^{-1})b(aa^{-1})b\cdots b(aa^{-1})ba^{-1} \\
&= ab^n a^{-1}.
\end{aligned}$$

3.4.31.   **Hint**.   Since $abab = (ab)^2 = e = a^2b^2 = aabb$, we know that $ba = ab$.

3.4.35.   **Hint**. $H_1 = \{\text{id}\}$, $H_2 = \{\text{id}, \rho_1, \rho_2\}$, $H_3 = \{\text{id}, \mu_1\}$, $H_4 = \{\text{id}, \mu_2\}$, $H_5 = \{\text{id}, \mu_3\}$, S_3.

3.4.41.   **Hint**.   The identity of G is $1 = 1 + 0\sqrt{2}$. Since $(a + b\sqrt{2})(c + d\sqrt{2}) = (ac + 2bd) + (ad + bc)\sqrt{2}$, G is closed under multiplication. Finally, $(a + b\sqrt{2})^{-1} = a/(a^2 - 2b^2) - b\sqrt{2}/(a^2 - 2b^2)$.

3.4.46.   **Hint**.   Look at S_3.

3.4.49.   **Hint**. $ba = a^4b = a^3ab = ab$

4 Cyclic Groups

4.4 Exercises

4.4.1.   **Hint**.   (a) False; (c) false; (e) true.

4.4.2.   **Hint**.   (a) 12; (c) infinite; (e) 10.

4.4.3.   **Hint**.   (a) $7\mathbb{Z} = \{\ldots, -7, 0, 7, 14, \ldots\}$; (b) $\{0, 3, 6, 9, 12, 15, 18, 21\}$; (c) $\{0\}$, $\{0, 6\}$, $\{0, 4, 8\}$, $\{0, 3, 6, 9\}$, $\{0, 2, 4, 6, 8, 10\}$; (g) $\{1, 3, 7, 9\}$; (j) $\{1, -1, i, -i\}$.

4.4.4.   **Hint**.   (a)

$$\begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}, \begin{pmatrix} -1 & 0 \\ 0 & -1 \end{pmatrix}, \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix}, \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}.$$

(c)

$$\begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}, \begin{pmatrix} 1 & -1 \\ 1 & 0 \end{pmatrix}, \begin{pmatrix} -1 & 1 \\ -1 & 0 \end{pmatrix}, \begin{pmatrix} 0 & 1 \\ -1 & 1 \end{pmatrix}, \begin{pmatrix} 0 & -1 \\ 1 & -1 \end{pmatrix}, \begin{pmatrix} -1 & 0 \\ 0 & -1 \end{pmatrix}.$$

4.4.10.   **Hint**.   (a) 0; (b) 1, -1.

4.4.11.   **Hint**. $1, 2, 3, 4, 6, 8, 12, 24$.

4.4.15.   **Hint**.   (a) $-3 + 3i$; (c) $43 - 18i$; (e) i

4.4.16.   **Hint**.   (a) $\sqrt{3} + i$; (c) -3.

4.4.17. **Hint.** (a) $\sqrt{2}\,\text{cis}(7\pi/4)$; (c) $2\sqrt{2}\,\text{cis}(\pi/4)$; (e) $3\,\text{cis}(3\pi/2)$.

4.4.18. **Hint.** (a) $(1-i)/2$; (c) $16(i-\sqrt{3})$; (e) $-1/4$.

4.4.22. **Hint.** (a) 292; (c) 1523.

4.4.27. **Hint.** $|\langle g \rangle \cap \langle h \rangle| = 1$.

4.4.31. **Hint.** The identity element in any group has finite order. Let $g, h \in G$ have orders m and n, respectively. Since $(g^{-1})^m = e$ and $(gh)^{mn} = e$, the elements of finite order in G form a subgroup of G.

4.4.37. **Hint.** If g is an element distinct from the identity in G, g must generate G; otherwise, $\langle g \rangle$ is a nontrivial proper subgroup of G.

5 Permutation Groups

5.3 Exercises

5.3.1. **Hint.** (a) (12453); (c) (13)(25).

5.3.2. **Hint.** (a) (135)(24); (c) (14)(23); (e) (1324); (g) (134)(25); (n) (17352).

5.3.3. **Hint.** (a) (16)(15)(13)(14); (c) (16)(14)(12).

5.3.4. **Hint.** $(a_1, a_2, \ldots, a_n)^{-1} = (a_1, a_n, a_{n-1}, \ldots, a_2)$

5.3.5. **Hint.** (a) $\{(13), (13)(24), (132), (134), (1324), (1342)\}$ is not a subgroup.

5.3.8. **Hint.** (12345)(678).

5.3.11. **Hint.** Permutations of the form

$$(1), (a_1, a_2)(a_3, a_4), (a_1, a_2, a_3), (a_1, a_2, a_3, a_4, a_5)$$

are possible for A_5.

5.3.17. **Hint.** Calculate $(123)(12)$ and $(12)(123)$.

5.3.25. **Hint.** Consider the cases $(ab)(bc)$ and $(ab)(cd)$.

5.3.30. **Hint.** For (a), show that $\sigma\tau\sigma^{-1}(\sigma(a_i)) = \sigma(a_{i+1})$.

6 Cosets and Lagrange's Theorem

6.4 Exercises

6.4.1. **Hint.** The order of g and the order h must both divide the order of G.

6.4.2. **Hint.** The possible orders must divide 60.

6.4.3. **Hint.** This is true for every proper nontrivial subgroup.

6.4.4. **Hint.** False.

6.4.5. **Hint.** (a) $\langle 8 \rangle$, $1 + \langle 8 \rangle$, $2 + \langle 8 \rangle$, $3 + \langle 8 \rangle$, $4 + \langle 8 \rangle$, $5 + \langle 8 \rangle$, $6 + \langle 8 \rangle$, and $7 + \langle 8 \rangle$; (c) $3\mathbb{Z}$, $1 + 3\mathbb{Z}$, and $2 + 3\mathbb{Z}$.

6.4.7. **Hint.** $4^{\phi(15)} \equiv 4^8 \equiv 1 \pmod{15}$.

6.4.12. **Hint.** Let $g_1 \in gH$. Show that $g_1 \in Hg$ and thus $gH \subset Hg$.

6.4.19. **Hint.** Show that $g(H \cap K) = gH \cap gK$.

6.4.22. **Hint.** If $\gcd(m, n) = 1$, then $\phi(mn) = \phi(m)\phi(n)$ (Exercise 2.3.26 in Chapter 2).

7 Introduction to Cryptography

7.3 Exercises

7.3.1.   **Hint**.   LAORYHAPDWK

7.3.3.   **Hint**.   Hint: V = E, E = X (also used for spaces and punctuation), K = R.

7.3.4.   **Hint**. $26! - 1$

7.3.7.   **Hint**.   (a) 2791; (c) 11213525032442.

7.3.9.   **Hint**.   (a) 31 (c) 14.

7.3.10.   **Hint**.   (a) $n = 11 \cdot 41$; (c) $n = 8779 \cdot 4327$.

8 Algebraic Coding Theory

8.5 Exercises

8.5.2.   **Hint**.   This cannot be a group code since $(0000) \notin C$.

8.5.3.   **Hint**.   (a) 2; (c) 2.

8.5.4.   **Hint**.   (a) 3; (c) 4.

8.5.6.   **Hint**.   (a) $d_{\min} = 2$; (c) $d_{\min} = 1$.

8.5.7.   **Hint**.

  (a) $(00000), (00101), (10011), (10110)$

$$G = \begin{pmatrix} 0 & 1 \\ 0 & 0 \\ 1 & 0 \\ 0 & 1 \\ 1 & 1 \end{pmatrix}$$

  (b) $(000000), (010111), (101101), (111010)$

$$G = \begin{pmatrix} 1 & 0 \\ 0 & 1 \\ 1 & 0 \\ 1 & 1 \\ 0 & 1 \\ 1 & 1 \end{pmatrix}$$

8.5.9.   **Hint**.   Multiple errors occur in one of the received words.

8.5.11.   **Hint**.   (a) A canonical parity-check matrix with standard generator matrix

$$G = \begin{pmatrix} 1 \\ 1 \\ 0 \\ 0 \\ 1 \end{pmatrix}.$$

(c) A canonical parity-check matrix with standard generator matrix

$$G = \begin{pmatrix} 1 & 0 \\ 0 & 1 \\ 1 & 1 \\ 1 & 0 \end{pmatrix}.$$

8.5.12. **Hint.** (a) All possible syndromes occur.

8.5.15. **Hint.** (a) C, $(10000) + C$, $(01000) + C$, $(00100) + C$, $(00010) + C$, $(11000) + C$, $(01100) + C$, $(01010) + C$. A decoding table does not exist for C since this is only a single error-detecting code.

8.5.19. **Hint.** Let $\mathbf{x} \in C$ have odd weight and define a map from the set of odd codewords to the set of even codewords by $\mathbf{y} \mapsto \mathbf{x} + \mathbf{y}$. Show that this map is a bijection.

8.5.23. **Hint.** For 20 information positions, at least 6 check bits are needed to ensure an error-correcting code.

9 Isomorphisms

9.3 Exercises

9.3.1. **Hint.** Every infinite cyclic group is isomorphic to $\mathbb{Z}$ by Theorem 9.7.

9.3.2. **Hint.** Define $\phi : \mathbb{C}^* \to GL_2(\mathbb{R})$ by

$$\phi(a + bi) = \begin{pmatrix} a & b \\ -b & a \end{pmatrix}.$$

9.3.3. **Hint.** False.

9.3.6. **Hint.** Define a map from $\mathbb{Z}_n$ into the nth roots of unity by $k \mapsto \operatorname{cis}(2k\pi/n)$.

9.3.8. **Hint.** Assume that $\mathbb{Q}$ is cyclic and try to find a generator.

9.3.11. **Hint.** There are two nonabelian and three abelian groups that are not isomorphic.

9.3.16. **Hint.** (a) 12; (c) 5.

9.3.19. **Hint.** Draw the picture.

9.3.20. **Hint.** True.

9.3.25. **Hint.** True.

9.3.27. **Hint.** Let a be a generator for G. If $\phi : G \to H$ is an isomorphism, show that $\phi(a)$ is a generator for H.

9.3.38. **Hint.** Any automorphism of $\mathbb{Z}_6$ must send 1 to another generator of $\mathbb{Z}_6$.

9.3.45. **Hint.** To show that ϕ is one-to-one, let $g_1 = h_1 k_1$ and $g_2 = h_2 k_2$ and consider $\phi(g_1) = \phi(g_2)$.

10 Normal Subgroups and Factor Groups

10.3 Exercises

10.3.1. **Hint**. (a)

	A_4	$(12)A_4$
A_4	A_4	$(12)A_4$
$(12)A_4$	$(12)A_4$	A_4

(c) D_4 is not normal in S_4.

10.3.8. **Hint**. If $a \in G$ is a generator for G, then aH is a generator for G/H.

10.3.11. **Hint**. For any $g \in G$, show that the map $i_g : G \to G$ defined by $i_g : x \mapsto gxg^{-1}$ is an isomorphism of G with itself. Then consider $i_g(H)$.

10.3.12. **Hint**. Suppose that $\langle g \rangle$ is normal in G and let y be an arbitrary element of G. If $x \in C(g)$, we must show that yxy^{-1} is also in $C(g)$. Show that $(yxy^{-1})g = g(yxy^{-1})$.

10.3.14. **Hint**. (a) Let $g \in G$ and $h \in G'$. If $h = aba^{-1}b^{-1}$, then

$$ghg^{-1} = gaba^{-1}b^{-1}g^{-1}$$
$$= (gag^{-1})(gbg^{-1})(ga^{-1}g^{-1})(gb^{-1}g^{-1})$$
$$= (gag^{-1})(gbg^{-1})(gag^{-1})^{-1}(gbg^{-1})^{-1}.$$

We also need to show that if $h = h_1 \cdots h_n$ with $h_i = a_i b_i a_i^{-1} b_i^{-1}$, then ghg^{-1} is a product of elements of the same type. However, $ghg^{-1} = gh_1 \cdots h_n g^{-1} = (gh_1 g^{-1})(gh_2 g^{-1}) \cdots (gh_n g^{-1})$.

11 Homomorphisms

11.3 Exercises

11.3.2. **Hint**. (a) is a homomorphism with kernel $\{1\}$; (c) is not a homomorphism.

11.3.4. **Hint**. Since $\phi(m + n) = 7(m + n) = 7m + 7n = \phi(m) + \phi(n)$, ϕ is a homomorphism.

11.3.5. **Hint**. For any homomorphism $\phi : \mathbb{Z}_{24} \to \mathbb{Z}_{18}$, the kernel of ϕ must be a subgroup of $\mathbb{Z}_{24}$ and the image of ϕ must be a subgroup of $\mathbb{Z}_{18}$. Now use the fact that a generator must map to a generator.

11.3.9. **Hint**. Let $a, b \in G$. Then $\phi(a)\phi(b) = \phi(ab) = \phi(ba) = \phi(b)\phi(a)$.

11.3.17. **Hint**. Find a counterexample.

12 Matrix Groups and Symmetry

12.3 Exercises

12.3.1. **Hint**.

$$\frac{1}{2} \left[\|\mathbf{x} + \mathbf{y}\|^2 + \|\mathbf{x}\|^2 - \|\mathbf{y}\|^2 \right] = \frac{1}{2} \left[\langle x + y, x + y \rangle - \|\mathbf{x}\|^2 - \|\mathbf{y}\|^2 \right]$$
$$= \frac{1}{2} \left[\|\mathbf{x}\|^2 + 2\langle x, y \rangle + \|\mathbf{y}\|^2 - \|\mathbf{x}\|^2 - \|\mathbf{y}\|^2 \right]$$
$$= \langle \mathbf{x}, \mathbf{y} \rangle.$$

12.3.3.   **Hint**.   (a) is in $SO(2)$; (c) is not in $O(3)$.

12.3.5.   **Hint**.   (a) $\langle \mathbf{x}, \mathbf{y} \rangle = \langle \mathbf{y}, \mathbf{x} \rangle$.

12.3.7.   **Hint**.   Use the unimodular matrix

$$\begin{pmatrix} 5 & 2 \\ 2 & 1 \end{pmatrix}.$$

12.3.10.   **Hint**.   Show that the kernel of the map $\det : O(n) \to \mathbb{R}^*$ is $SO(n)$.

12.3.13.   **Hint**.   True.

12.3.17.   **Hint**. $p6m$

13 The Structure of Groups

13.3 Exercises

13.3.1.   **Hint**.   There are three possible groups.

13.3.4.   **Hint**.   (a) $\{0\} \subset \langle 6 \rangle \subset \langle 3 \rangle \subset \mathbb{Z}_{12}$; (e) $\{(1)\} \times \{0\} \subset \{(1), (123), (132)\} \times \{0\} \subset S_3 \times \{0\} \subset S_3 \times \langle 2 \rangle \subset S_3 \times \mathbb{Z}_4$.

13.3.7.   **Hint**.   Use the Fundamental Theorem of Finitely Generated Abelian Groups.

13.3.12.   **Hint**.   If N and G/N are solvable, then they have solvable series

$$N = N_n \supset N_{n-1} \supset \cdots \supset N_1 \supset N_0 = \{e\}$$
$$G/N = G_n/N \supset G_{n-1}/N \supset \cdots G_1/N \supset G_0/N = \{N\}.$$

13.3.16.   **Hint**.   Use the fact that D_n has a cyclic subgroup of index 2.

13.3.21.   **Hint**. G/G' is abelian.

14 Group Actions

14.4 Exercises

14.4.1.   **Hint**.   Example 14.1: 0, $\mathbb{R}^2 \setminus \{0\}$. Example 14.2: $X = \{1, 2, 3, 4\}$.

14.4.2.   **Hint**.   (a) $X_{(1)} = \{1, 2, 3\}$, $X_{(12)} = \{3\}$, $X_{(13)} = \{2\}$, $X_{(23)} = \{1\}$, $X_{(123)} = X_{(132)} = \emptyset$. $G_1 = \{(1), (23)\}$, $G_2 = \{(1), (13)\}$, $G_3 = \{(1), (12)\}$.

14.4.3.   **Hint**.   (a) $\mathcal{O}_1 = \mathcal{O}_2 = \mathcal{O}_3 = \{1, 2, 3\}$.

14.4.6.   **Hint**.   The conjugacy classes for S_4 are

$$\mathcal{O}_{(1)} = \{(1)\},$$
$$\mathcal{O}_{(12)} = \{(12), (13), (14), (23), (24), (34)\},$$
$$\mathcal{O}_{(12)(34)} = \{(12)(34), (13)(24), (14)(23)\},$$
$$\mathcal{O}_{(123)} = \{(123), (132), (124), (142), (134), (143), (234), (243)\},$$
$$\mathcal{O}_{(1234)} = \{(1234), (1243), (1324), (1342), (1423), (1432)\}.$$

The class equation is $1 + 3 + 6 + 6 + 8 = 24$.

14.4.8.   **Hint**. $(3^4 + 3^1 + 3^2 + 3^1 + 3^2 + 3^2 + 3^3 + 3^3)/8 = 21$.

14.4.11.   **Hint**.   The group of rigid motions of the cube can be described by the allowable permutations of the six faces and is isomorphic to S_4. There are the identity cycle, 6 permutations with the structure $(abcd)$ that correspond to the quarter turns, 3 permutations with the structure $(ab)(cd)$ that correspond to the half turns, 6 permutations with the structure $(ab)(cd)(ef)$ that correspond to rotating the cube about the centers of opposite edges, and 8 permutations with the structure $(abc)(def)$ that correspond to rotating the cube about opposite vertices.

14.4.15.   **Hint**. $(1 \cdot 2^6 + 3 \cdot 2^4 + 4 \cdot 2^3 + 2 \cdot 2^2 + 2 \cdot 2^1)/12 = 13$.

14.4.17.   **Hint**. $(1 \cdot 2^8 + 3 \cdot 2^6 + 2 \cdot 2^4)/6 = 80$.

14.4.22.   **Hint**.   Use the fact that $x \in gC(a)g^{-1}$ if and only if $g^{-1}xg \in C(a)$.

15 The Sylow Theorems

15.3 Exercises

15.3.1.   **Hint**.   If $|G| = 18 = 2 \cdot 3^2$, then the order of a Sylow 2-subgroup is 2, and the order of a Sylow 3-subgroup is 9.

15.3.2.   **Hint**.   The four Sylow 3-subgroups of S_4 are $P_1 = \{(1), (123), (132)\}$, $P_2 = \{(1), (124), (142)\}$, $P_3 = \{(1), (134), (143)\}$, $P_4 = \{(1), (234), (243)\}$.

15.3.5.   **Hint**.   Since $|G| = 96 = 2^5 \cdot 3$, G has either one or three Sylow 2-subgroups by the Third Sylow Theorem. If there is only one subgroup, we are done. If there are three Sylow 2-subgroups, let H and K be two of them. Therefore, $|H \cap K| \geq 16$; otherwise, HK would have $(32 \cdot 32)/8 = 128$ elements, which is impossible. Thus, $H \cap K$ is normal in both H and K since it has index 2 in both groups.

15.3.8.   **Hint**.   Show that G has a normal Sylow p-subgroup of order p^2 and a normal Sylow q-subgroup of order q^2.

15.3.10.   **Hint**.   False.

15.3.17.   **Hint**.   If G is abelian, then G is cyclic, since $|G| = 3 \cdot 5 \cdot 17$. Now look at Example 15.14.

15.3.23.   **Hint**.   Define a mapping between the right cosets of $N(H)$ in G and the conjugates of H in G by $N(H)g \mapsto g^{-1}Hg$. Prove that this map is a bijection.

15.3.26.   **Hint**.   Let $aG', bG' \in G/G'$. Then $(aG')(bG') = abG' = ab(b^{-1}a^{-1}ba)G' = (abb^{-1}a^{-1})baG' = baG'$.

16 Rings

16.6 Exercises

16.6.1.   **Hint**.   (a) $7\mathbb{Z}$ is a ring but not a field; (c) $\mathbb{Q}(\sqrt{2})$ is a field; (f) R is not a ring.

16.6.3.   **Hint**.   (a) $\{1, 3, 7, 9\}$; (c) $\{1, 2, 3, 4, 5, 6\}$; (e)

$$\left\{ \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}, \begin{pmatrix} 1 & 1 \\ 0 & 1 \end{pmatrix}, \begin{pmatrix} 1 & 0 \\ 1 & 1 \end{pmatrix}, \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix}, \begin{pmatrix} 1 & 1 \\ 1 & 0 \end{pmatrix}, \begin{pmatrix} 0 & 1 \\ 1 & 1 \end{pmatrix}, \right\}.$$

16.6.4.   **Hint**.   (a) $\{0\}$, $\{0, 9\}$, $\{0, 6, 12\}$, $\{0, 3, 6, 9, 12, 15\}$, $\{0, 2, 4, 6, 8, 10, 12, 14, 16\}$; (c) there are no nontrivial ideals.

16.6.7.   **Hint**.   Assume there is an isomorphism $\phi : \mathbb{C} \to \mathbb{R}$ with $\phi(i) = a$.

16.6.8. **Hint.** False. Assume there is an isomorphism $\phi : \mathbb{Q}(\sqrt{2}) \to \mathbb{Q}(\sqrt{3})$ such that $\phi(\sqrt{2}) = a$.

16.6.13. **Hint.** (a) $x \equiv 17 \pmod{55}$; (c) $x \equiv 214 \pmod{2772}$.

16.6.16. **Hint.** If $I \neq \{0\}$, show that $1 \in I$.

16.6.18. **Hint.** (a) $\phi(a)\phi(b) = \phi(ab) = \phi(ba) = \phi(b)\phi(a)$.

16.6.26. **Hint.** Let $a \in R$ with $a \neq 0$. Then the principal ideal generated by a is R. Thus, there exists a $b \in R$ such that $ab = 1$.

16.6.28. **Hint.** Compute $(a + b)^2$ and $(-ab)^2$.

16.6.34. **Hint.** Let $a/b, c/d \in \mathbb{Z}_{(p)}$. Then $a/b + c/d = (ad + bc)/bd$ and $(a/b) \cdot (c/d) = (ac)/(bd)$ are both in $\mathbb{Z}_{(p)}$, since $\gcd(bd, p) = 1$.

16.6.38. **Hint.** Suppose that $x^2 = x$ and $x \neq 0$. Since R is an integral domain, $x = 1$. To find a nontrivial idempotent, look in $\mathbb{M}_2(\mathbb{R})$.

17 Polynomials

17.4 Exercises

17.4.2. **Hint.** (a) $9x^2 + 2x + 5$; (b) $8x^4 + 7x^3 + 2x^2 + 7x$.

17.4.3. **Hint.** (a) $5x^3 + 6x^2 - 3x + 4 = (5x^2 + 2x + 1)(x - 2) + 6$; (c) $4x^5 - x^3 + x^2 + 4 = (4x^2 + 4)(x^3 + 3) + 4x^2 + 2$.

17.4.5. **Hint.** (a) No zeros in $\mathbb{Z}_{12}$; (c) 3, 4.

17.4.7. **Hint.** Look at $(2x + 1)$.

17.4.8. **Hint.** (a) Reducible; (c) irreducible.

17.4.10. **Hint.** One factorization is $x^2 + x + 8 = (x + 2)(x + 9)$.

17.4.13. **Hint.** The integers $\mathbb{Z}$ do not form a field.

17.4.14. **Hint.** False.

17.4.16. **Hint.** Let $\phi : R \to S$ be an isomorphism. Define $\overline{\phi} : R[x] \to S[x]$ by $\overline{\phi}(a_0 + a_1 x + \cdots + a_n x^n) = \phi(a_0) + \phi(a_1)x + \cdots + \phi(a_n)x^n$.

17.4.20. Cyclotomic Polynomials. **Hint.** The polynomial

$$\Phi_n(x) = \frac{x^n - 1}{x - 1} = x^{n-1} + x^{n-2} + \cdots + x + 1$$

is called the ***cyclotomic polynomial.*** Show that $\Phi_p(x)$ is irreducible over $\mathbb{Q}$ for any prime p.

17.4.26. **Hint.** Find a nontrivial proper ideal in $F[x]$.

18 Integral Domains

18.3 Exercises

18.3.1. **Hint.** Note that $z^{-1} = 1/(a + b\sqrt{3}\,i) = (a - b\sqrt{3}\,i)/(a^2 + 3b^2)$ is in $\mathbb{Z}[\sqrt{3}\,i]$ if and only if $a^2 + 3b^2 = 1$. The only integer solutions to the equation are $a = \pm 1, b = 0$.

18.3.2. **Hint.** (a) $5 = -i(1 + 2i)(2 + i)$; (c) $6 + 8i = -i(1 + i)^2(2 + i)^2$.

18.3.4. **Hint.** True.

18.3.9. **Hint.** Let $z = a + bi$ and $w = c + di \neq 0$ be in $\mathbb{Z}[i]$. Prove that $z/w \in \mathbb{Q}(i)$.

18.3.15.   **Hint**.   Let $a = ub$ with u a unit. Then $\nu(b) \le \nu(ub) \le \nu(a)$. Similarly, $\nu(a) \le \nu(b)$.

18.3.16.   **Hint**.   Show that 21 can be factored in two different ways.

19 Lattices and Boolean Algebras

19.4 Exercises

19.4.2.   **Hint**.
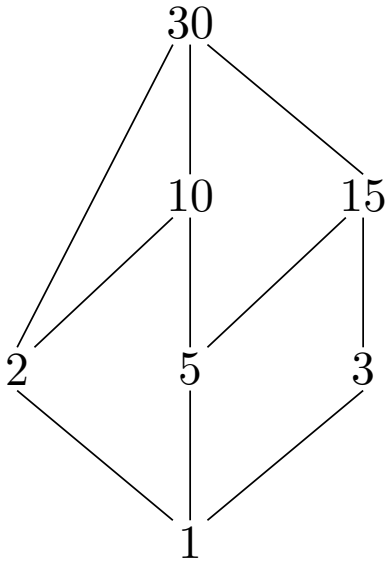
19.4.5.   **Hint**.   False.

19.4.6.   **Hint**.   (a) $(a \vee b \vee a') \wedge a$
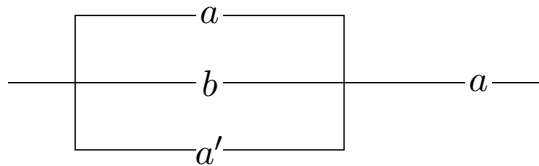
(c) $a \vee (a \wedge b)$
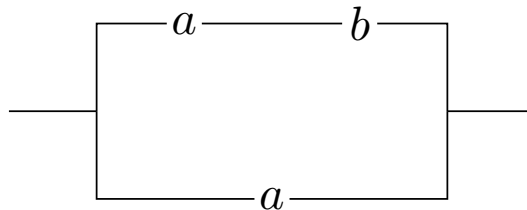
19.4.8.   **Hint**.   Not equivalent.

19.4.10.   **Hint**.   (a) $a' \wedge [(a \wedge b') \vee b] = a \wedge (a \vee b)$.

19.4.14. **Hint.** Let I, J be ideals in R. We need to show that $I + J = \{r + s : r \in I \text{ and } s \in J\}$ is the smallest ideal in R containing both I and J. If $r_1, r_2 \in I$ and $s_1, s_2 \in J$, then $(r_1 + s_1) + (r_2 + s_2) = (r_1 + r_2) + (s_1 + s_2)$ is in $I + J$. For $a \in R$, $a(r_1 + s_1) = ar_1 + as_1 \in I + J$; hence, $I + J$ is an ideal in R.

19.4.18. **Hint.** (a) No.

19.4.20. **Hint.** ($\Rightarrow$). $a = b \Rightarrow (a \wedge b') \vee (a' \wedge b) = (a \wedge a') \vee (a' \wedge a) = O \vee O = O$. ($\Leftarrow$). $(a \wedge b') \vee (a' \wedge b) = O \Rightarrow a \vee b = (a \vee a) \vee b = a \vee (a \vee b) = a \vee [I \wedge (a \vee b)] = a \vee [(a \vee a') \wedge (a \vee b)] = [a \vee (a \wedge b')] \vee [a \vee (a' \wedge b)] = a \vee [(a \wedge b') \vee (a' \wedge b)] = a \vee 0 = a$. A symmetric argument shows that $a \vee b = b$.

20 Vector Spaces

20.4 Exercises

20.4.3. **Hint.** $\mathbb{Q}(\sqrt{2}, \sqrt{3})$ has basis $\{1, \sqrt{2}, \sqrt{3}, \sqrt{6}\}$ over $\mathbb{Q}$.

20.4.5. **Hint.** The set $\{1, x, x^2, \ldots, x^{n-1}\}$ is a basis for P_n.

20.4.7. **Hint.** (a) Subspace of dimension 2 with basis $\{(1, 0, -3), (0, 1, 2)\}$; (d) not a subspace

20.4.10. **Hint.** Since $0 = \alpha 0 = \alpha(-v + v) = \alpha(-v) + \alpha v$, it follows that $-\alpha v = \alpha(-v)$.

20.4.12. **Hint.** Let $v_0 = 0, v_1, \ldots, v_n \in V$ and $\alpha_0 \neq 0, \alpha_1, \ldots, \alpha_n \in F$. Then $\alpha_0 v_0 + \cdots + \alpha_n v_n = 0$.

20.4.15. **Linear Transformations.** **Hint.** (a) Let $u, v \in \ker(T)$ and $\alpha \in F$. Then

$$T(u + v) = T(u) + T(v) = 0$$
$$T(\alpha v) = \alpha T(v) = \alpha 0 = 0.$$

Hence, $u + v, \alpha v \in \ker(T)$, and $\ker(T)$ is a subspace of V.

(c) The statement that $T(u) = T(v)$ is equivalent to $T(u - v) = T(u) - T(v) = 0$, which is true if and only if $u - v = 0$ or $u = v$.

20.4.17. **Direct Sums.** **Hint.** (a) Let $u, u' \in U$ and $v, v' \in V$. Then

$$(u + v) + (u' + v') = (u + u') + (v + v') \in U + V$$
$$\alpha(u + v) = \alpha u + \alpha v \in U + V.$$

21 Fields

21.4 Exercises

21.4.1. **Hint.** (a) $x^4 - (2/3)x^2 - 62/9$; (c) $x^4 - 2x^2 + 25$.

21.4.2. **Hint.** (a) $\{1, \sqrt{2}, \sqrt{3}, \sqrt{6}\}$; (c) $\{1, i, \sqrt{2}, \sqrt{2}\,i\}$; (e) $\{1, 2^{1/6}, 2^{1/3}, 2^{1/2}, 2^{2/3}, 2^{5/6}\}$.

21.4.3. **Hint.** (a) $\mathbb{Q}(\sqrt{3}, \sqrt{7})$.

21.4.5. **Hint.** Use the fact that the elements of $\mathbb{Z}_2[x]/\langle x^3 + x + 1\rangle$ are 0, 1, α, $1 + \alpha$, α^2, $1 + \alpha^2$, $\alpha + \alpha^2$, $1 + \alpha + \alpha^2$ and the fact that $\alpha^3 + \alpha + 1 = 0$.

21.4.8. **Hint.** False.

21.4.14. **Hint.** Suppose that E is algebraic over F and K is algebraic over E. Let $\alpha \in K$. It suffices to show that α is algebraic over some finite extension of F. Since α is

algebraic over E, it must be the zero of some polynomial $p(x) = \beta_0 + \beta_1 x + \cdots + \beta_n x^n$ in $E[x]$. Hence α is algebraic over $F(\beta_0, \ldots, \beta_n)$.

21.4.22.   **Hint.**   Since $\{1, \sqrt{3}, \sqrt{7}, \sqrt{21}\,\}$ is a basis for $\mathbb{Q}(\sqrt{3}, \sqrt{7}\,)$ over $\mathbb{Q}$, $\mathbb{Q}(\sqrt{3}, \sqrt{7}\,) \supset \mathbb{Q}(\sqrt{3} + \sqrt{7}\,)$. Since $[\mathbb{Q}(\sqrt{3}, \sqrt{7}\,) : \mathbb{Q}] = 4$, $[\mathbb{Q}(\sqrt{3} + \sqrt{7}\,) : \mathbb{Q}] = 2$ or 4. Since the degree of the minimal polynomial of $\sqrt{3} + \sqrt{7}$ is 4, $\mathbb{Q}(\sqrt{3}, \sqrt{7}\,) = \mathbb{Q}(\sqrt{3} + \sqrt{7}\,)$.

21.4.27.   **Hint.**   Let $\beta \in F(\alpha)$ not in F. Then $\beta = p(\alpha)/q(\alpha)$, where p and q are polynomials in α with $q(\alpha) \neq 0$ and coefficients in F. If β is algebraic over F, then there exists a polynomial $f(x) \in F[x]$ such that $f(\beta) = 0$. Let $f(x) = a_0 + a_1 x + \cdots + a_n x^n$. Then

$$0 = f(\beta) = f\left(\frac{p(\alpha)}{q(\alpha)}\right) = a_0 + a_1 \left(\frac{p(\alpha)}{q(\alpha)}\right) + \cdots + a_n \left(\frac{p(\alpha)}{q(\alpha)}\right)^n.$$

Now multiply both sides by $q(\alpha)^n$ to show that there is a polynomial in $F[x]$ that has α as a zero.

21.4.28.   **Hint.**   See the comments following Theorem 21.13.

22 Finite Fields

22.3 Exercises

22.3.1.   **Hint.**   Make sure that you have a field extension.

22.3.4.   **Hint.**   There are eight elements in $\mathbb{Z}_2(\alpha)$. Exhibit two more zeros of $x^3 + x^2 + 1$ other than α in these eight elements.

22.3.5.   **Hint.**   Find an irreducible polynomial $p(x)$ in $\mathbb{Z}_3[x]$ of degree 3 and show that $\mathbb{Z}_3[x]/\langle p(x)\rangle$ has 27 elements.

22.3.7.   **Hint.**   (a) $x^5 - 1 = (x + 1)(x^4 + x^3 + x^2 + x + 1)$; (c) $x^9 - 1 = (x + 1)(x^2 + x + 1)(x^6 + x^3 + 1)$.

22.3.8.   **Hint.**   True.

22.3.11.   **Hint.**   (a) Use the fact that $x^7 - 1 = (x + 1)(x^3 + x + 1)(x^3 + x^2 + 1)$.

22.3.12.   **Hint.**   False.

22.3.17.   **Hint.**   If $p(x) \in F[x]$, then $p(x) \in E[x]$.

22.3.18.   **Hint.**   Since α is algebraic over F of degree n, we can write any element $\beta \in F(\alpha)$ uniquely as $\beta = a_0 + a_1 \alpha + \cdots + a_{n-1} \alpha^{n-1}$ with $a_i \in F$. There are q^n possible n-tuples $(a_0, a_1, \ldots, a_{n-1})$.

22.3.24. Wilson's Theorem.   **Hint.**   Factor $x^{p-1} - 1$ over $\mathbb{Z}_p$.

23 Galois Theory

23.4 Exercises

23.4.1.   **Hint.**   (a) $\mathbb{Z}_2$; (c) $\mathbb{Z}_2 \times \mathbb{Z}_2 \times \mathbb{Z}_2$.

23.4.2.   **Hint.**   (a) Separable over $\mathbb{Q}$ since $x^3 + 2x^2 - x - 2 = (x - 1)(x + 1)(x + 2)$; (c) not separable over $\mathbb{Z}_3$ since $x^4 + x^2 + 1 = (x + 1)^2(x + 2)^2$.

23.4.3.   **Hint.**   If

$$[\mathrm{GF}(729) : \mathrm{GF}(9)] = [\mathrm{GF}(729) : \mathrm{GF}(3)]/[\mathrm{GF}(9) : \mathrm{GF}(3)] = 6/2 = 3,$$

then $G(\mathrm{GF}(729)/\mathrm{GF}(9)) \cong \mathbb{Z}_3$. A generator for $G(\mathrm{GF}(729)/\mathrm{GF}(9))$ is σ, where $\sigma_{3^6}(\alpha) = \alpha^{3^6} = \alpha^{729}$ for $\alpha \in \mathrm{GF}(729)$.

23.4.4. **Hint**. (a) S_5; (c) S_3; (g) see Example 23.10.

23.4.5. **Hint**. (a) $\mathbb{Q}(i)$

23.4.7. **Hint**. Let E be the splitting field of a cubic polynomial in $F[x]$. Show that $[E : F]$ is less than or equal to 6 and is divisible by 3. Since $G(E/F)$ is a subgroup of S_3 whose order is divisible by 3, conclude that this group must be isomorphic to $\mathbb{Z}_3$ or S_3.

23.4.9. **Hint**. G is a subgroup of S_n.

23.4.16. **Hint**. True.

23.4.20. **Hint**.

(a) Clearly $\omega, \omega^2, \ldots, \omega^{p-1}$ are distinct since $\omega \neq 1$ or 0. To show that ω^i is a zero of Φ_p, calculate $\Phi_p(\omega^i)$.

(b) The conjugates of ω are $\omega, \omega^2, \ldots, \omega^{p-1}$. Define a map $\phi_i : \mathbb{Q}(\omega) \to \mathbb{Q}(\omega^i)$ by

$$\phi_i(a_0 + a_1\omega + \cdots + a_{p-2}\omega^{p-2}) = a_0 + a_1\omega^i + \cdots + c_{p-2}(\omega^i)^{p-2},$$

where $a_i \in \mathbb{Q}$. Prove that ϕ_i is an isomorphism of fields. Show that ϕ_2 generates $G(\mathbb{Q}(\omega)/\mathbb{Q})$.

(c) Show that $\{\omega, \omega^2, \ldots, \omega^{p-1}\}$ is a basis for $\mathbb{Q}(\omega)$ over $\mathbb{Q}$, and consider which linear combinations of $\omega, \omega^2, \ldots, \omega^{p-1}$ are left fixed by all elements of $G(\mathbb{Q}(\omega)/\mathbb{Q})$.

$\boxed{C}$

Notation

The following table defines the notation used in this book. Page numbers or references refer to the first appearance of each symbol.

Symbol	Description	Page
$a \in A$	a is in the set A	3
$\mathbb{N}$	the natural numbers	4
$\mathbb{Z}$	the integers	4
$\mathbb{Q}$	the rational numbers	4
$\mathbb{R}$	the real numbers	4
$\mathbb{C}$	the complex numbers	4
$A \subset B$	A is a subset of B	4
$\emptyset$	the empty set	4
$A \cup B$	the union of sets A and B	4
$A \cap B$	the intersection of sets A and B	4
A'	complement of the set A	4
$A \setminus B$	difference between sets A and B	4
$A \times B$	Cartesian product of sets A and B	6
A^n	$A \times \cdots \times A$ (n times)	6
id	identity mapping	9
f^{-1}	inverse of the function f	9
$a \equiv b \pmod{n}$	a is congruent to b modulo n	12
$n!$	n factorial	23
$\binom{n}{k}$	binomial coefficient $n!/(k!(n-k)!)$	23
$a \mid b$	a divides b	25
$\gcd(a, b)$	greatest common divisor of a and b	25
$\mathcal{P}(X)$	power set of X	29
$\operatorname{lcm}(m, n)$	the least common multiple of m and n	30
$\mathbb{Z}_n$	the integers modulo n	36
$U(n)$	group of units in $\mathbb{Z}_n$	41
$\mathbb{M}_n(\mathbb{R})$	the $n \times n$ matrices with entries in $\mathbb{R}$	42
$\det A$	the determinant of A	42
$GL_n(\mathbb{R})$	the general linear group	42
Q_8	the group of quaternions	42
$\mathbb{C}^*$	the multiplicative group of complex numbers	42
$|G|$	the order of a group	43

(Continued on next page)

448

Symbol	Description	Page
$\mathbb{R}^*$	the multiplicative group of real numbers	45
$\mathbb{Q}^*$	the multiplicative group of rational numbers	45
$SL_n(\mathbb{R})$	the special linear group	45
$Z(G)$	the center of a group	49
$\langle a \rangle$	cyclic group generated by a	59
$|a|$	the order of an element a	59
$\operatorname{cis}\theta$	$\cos\theta + i\sin\theta$	63
$\mathbb{T}$	the circle group	64
S_n	the symmetric group on n letters	81
$(a_1, a_2, \ldots, a_k)$	cycle of length k	83
A_n	the alternating group on n letters	86
D_n	the dihedral group	87
$[G : H]$	index of a subgroup H in a group G	103
$\mathcal{L}_H$	the set of left cosets of a subgroup H in a group G	104
$\mathcal{R}_H$	the set of right cosets of a subgroup H in a group G	104
$a \nmid b$	a does not divide b	106
$d(\mathbf{x}, \mathbf{y})$	Hamming distance between $\mathbf{x}$ and $\mathbf{y}$	131
$d_{\min}$	the minimum distance of a code	131
$w(\mathbf{x})$	the weight of $\mathbf{x}$	131
$\mathbb{M}_{m \times n}(\mathbf{Z}_2)$	the set of $m \times n$ matrices with entries in $\mathbb{Z}_2$	135
$\operatorname{Null}(H)$	null space of a matrix H	135
δ_{ij}	Kronecker delta	139
$G \cong H$	G is isomorphic to a group H	153
$\operatorname{Aut}(G)$	automorphism group of a group G	162
i_g	$i_g(x) = gxg^{-1}$	162
$\operatorname{Inn}(G)$	inner automorphism group of a group G	162
ρ_g	right regular representation	162
G/N	factor group of G mod N	170
G'	commutator subgroup of G	175
$\ker \phi$	kernel of ϕ	182
(a_{ij})	matrix	195
$O(n)$	orthogonal group	197
$\|\mathbf{x}\|$	length of a vector $\mathbf{x}$	197
$SO(n)$	special orthogonal group	200
$E(n)$	Euclidean group	200
$\mathcal{O}_x$	orbit of x	223
X_g	fixed point set of g	223
G_x	isotropy subgroup of x	223
$N(H)$	normalizer of s subgroup H	241
$\mathbb{H}$	the ring of quaternions	258
$\mathbb{Z}[i]$	the Gaussian integers	259
$\operatorname{char} R$	characteristic of a ring R	260
$\mathbb{Z}_{(p)}$	ring of integers localized at p	272
$\deg f(x)$	degree of a polynomial	283

(Continued on next page)

Symbol	Description	Page
$R[x]$	ring of polynomials over a ring R	284
$R[x_1, x_2, \ldots, x_n]$	ring of polynomials in n indeterminants	286
ϕ_α	evaluation homomorphism at α	286
$\mathbb{Q}(x)$	field of rational functions over $\mathbb{Q}$	307
$\nu(a)$	Euclidean valuation of a	310
$F(x)$	field of rational functions in x	314
$F(x_1, \ldots, x_n)$	field of rational functions in $x_1, \ldots, x_n$	314
$a \preceq b$	a is less than b	320
$a \vee b$	join of a and b	322
$a \wedge b$	meet of a and b	322
I	largest element in a lattice	323
O	smallest element in a lattice	323
a'	complement of a in a lattice	323
$\dim V$	dimension of a vector space V	344
$U \oplus V$	direct sum of vector spaces U and V	346
$\mathrm{Hom}(V, W)$	set of all linear transformations from U into V	346
V^*	dual of a vector space V	346
$F(\alpha_1, \ldots, \alpha_n)$	smallest field containing F and $\alpha_1, \ldots, \alpha_n$	356
$[E : F]$	dimension of a field extension of E over F	359
$\mathrm{GF}(p^n)$	Galois field of order p^n	382
F^*	multiplicative group of a field F	383
$G(E/F)$	Galois group of E over F	398
$F_{\{\sigma_i\}}$	field fixed by the automorphism σ_i	401
F_G	field fixed by the automorphism group G	401
Δ^2	discriminant of a polynomial	413

Index

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