1. Introduction

The word “algebra” is derived from the title of a book - *Hisab al-jabr w'al-muqabala* - written by the Farsi scholar Abu Ja’far Muhammad ibn Musa Al-Khwarizmi (790 - 840). The word *al-jabr* itself comes from the root of *reunite* (in the sense of completing or putting together) and refers to one of the methods of solving quadratic equations (by completing the square) described in the book. The book can be considered as the first treatise on algebra. The word *algorithm* is in fact derived from the name *Al-Khwarizmi*. The book was very much concerned with recipes for known practical problems; Al-Khwarizmi intended to teach (in his own words) “… what is easiest and most useful in arithmetic, such as men constantly require in cases of inheritance, legacies, partition, lawsuits, and trade, and in all their dealings with one another, or where the measuring of lands, the digging of canals, geometrical computations, and other objects of various sorts and kinds are concerned.”

MATH 235 is a first course in Algebra. Little is assumed in way of background. Though the course is self-contained, it puts some of the responsibility of digesting and exploring the material on the student, as is normal in university studies. You’ll soon realize that we are also learning a new language in this course and a new attitude towards mathematics. The language is the language of modern mathematics; it is very formal, precise and concise. One of the challenges of the course is digesting and memorizing the new concepts and definitions. The new attitude is an attitude where any assumptions one is making while making an argument have to be justified, or at least clearly stated as a postulate, and from there on one proceeds in a logical and clear manner towards the conclusion. This is called “proof” and one of the main challenges in the course is to understand what constitutes a good proof and to be able to write proofs yourself. A further challenge for most students is that the key ideas we learn in this course are very abstract, bordering on philosophy and art, but yet they are truly scientific in their precision. You should expect not to understand everything right away; you should expect to need time to reflect on the meaning of the new ideas and concepts we introduce. Here are some pointers as to how to cope with the challenges of this course:

- Read the class notes and the text book over and over again. Try and give yourself examples of the theorems and propositions and try and provide counterexamples when some of the hypotheses are dropped.
- Do lots and lots of exercises. The more, the better.
- Explain to your friends, and possibly to your family, the material of the course. Work together with your class mates on assignments, but write your own solution by yourself in the end; try to understand different solutions to assignments and try and find flaws in your class mates solutions.
- Use the instructor’s and the TA’s office hours, as well as the math help center, to quickly close any gap and clarify any point you’re not sure about.

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1Cited from http://www-groups.dcs.st-and.ac.uk/ history/Biographies/Al-Khwarizmi.html.
So what is this course really about?
We are going to start by learning some of the notation and language of mathematics. We are
going to discuss sets and functions and various properties and operations one can perform on
those. We are going to talk about proofs and some techniques of proof, such as “induction”
and “proving the counter-positive”. We are going to discuss different structures in which one
can do arithmetic, such as rational, real and complex numbers, polynomial rings and yet more
abstract systems called rings and fields. We are going to see unifying patterns of such systems
and some of their applications. A finite field, perhaps a seemingly esoteric beast initially, is a key
notion in modern days computer science. It is a concept absolutely essential and fundamental
for cryptographic schemes as well as data management.

We are then going to do some really abstract algebra for a while, learning about rings and
homomorphisms and ideals. Our motivation is mostly to know how to construct a finite field and
work with it in practice. We also lay the basis for further study in the following algebra courses.

The final section of the course deals with groups and group actions on sets. After the previous
section on rings and fields we’ll feel more comfortable with the abstract notions of group theory.
Furthermore, there are going to be plenty of concrete examples in this section and applications
to problems in combinatorics and to the study of symmetries. Here is one concrete example:
Imagine that a jewelry company wants, as publicity, to display all necklaces one can make using
10 diamonds and 10 rubies, perhaps under the slogan “to each its own”. In each necklace 10
diamonds and 10 rubies are to be used. The necklace itself is just round, with no hanging or
protruding parts. Thus, we can provide an example of such a necklace as
\[ \text{D D R R D R D R R D D D R D R D R} \]
(where the last R is adjacent to the first D). Now, when we consider a particular design such as
above, we do want to identify it with the following design
\[ \text{D R R D R D R D R R D D D D R R D R D R} \]
(which we’ve put the first D at the last spot), because this is just the same pattern; if the necklace is
put on a table, it is just rotating it a bit, or, alternately, looking at it from a slightly different
angle. Also, note that the pattern
\[ \text{D D R R D R D R R D D D D R R D R D R} \]
is identified with
\[ \text{R D R D R D D D D R R D R D R D D} \]
which corresponds to flipping over the necklace. Now the question is how many rubies and
diamonds we need to purchase in order to make all the different designs? It turns out that this
can be approached using the theory of group actions on sets and a general formula we’ll develop
can be applied here. Turns out there are 4752 such different designs; that will require 47520
diamonds and 47520 rubies. Perhaps the idea should be reconsidered ;-)
Part 1. Some Language and Notation of Mathematics

2. Sets

A set is a collection of elements. The notion of a set is logically not quite defined (what’s a “collection”? an “element”?) but, hopefully, it makes sense to us. What we have is the ability to say whether an element is a member of a set or not. Thus, in a sense, a set is a property, and its elements are the objects having that property.

There are various ways to define sets:

(1) By writing it down:

\[ S = \{1, 3, 5\}. \]

The set is named \( S \) and its elements are 1, 3, 5. The use of curly brackets is mandatory!

Another example is

\[ T = \{2, 3, \text{Jim’s football}\}. \]

This is a set whose elements are the numbers 2, 3 and Jim’s football. It is assumed here that “Jim” refers to one particular individual.

A set can also be given as all objects with a certain property:

\[ S_1 = \{\text{all beluga whales}\}. \]

Another example is

\[ T_5 = \{n : n \text{ is an odd integer, } n^3 = n\}. \]

The colon means that the part that follows is the list of properties \( n \) must satisfy, i.e. the colon is shorthand to “such as”. Note that this set is equal to the set

\[ U^+ = \{n : n^2 = 1\}. \]

Our eccentric notation \( T_5, S_1, U^+ \) is just to make a point that a set can be denoted in many ways.

(2) Sometimes we write a set where the description of its elements is implicit, to be understood by the reader. For example:

\[ \mathbb{N} = \{0, 1, 2, 3, \ldots\}, \quad \mathbb{Z} = \{\ldots, -2, -1, 0, 1, 2, \ldots\}, \]

and

\[ \mathbb{Q} = \left\{ \frac{a}{b} : a, b \in \mathbb{Z}, b \neq 0 \right\}. \]

Thus \( \mathbb{N} \) is the set of natural numbers, \( \mathbb{Z} \) is the set of integers and \( \mathbb{Q} \) the set of rational numbers. The use of the letters \( \mathbb{N}, \mathbb{Z}, \mathbb{Q} \) (and \( \mathbb{R}, \mathbb{C} \) below) is standard. Other standard notation is

\[ \mathbb{R} = \text{the set of real numbers (} = \text{points on the line)}, \]
and the complex numbers
\[ 
\mathbb{C} = \{a + bi : a, b \in \mathbb{R}\}. 
\]

Here \( i \) is the imaginary number satisfying \( i^2 = -1 \) (we’ll come back to that in \( \S 6 \)). Note that we sneaked in new notation. If \( A \) is a set, the notation \( x \in A \) means \( x \) is an element (a member) of \( A \), while \( x \notin A \) means that \( x \) is not an element of \( A \). Thus, \( \mathbb{C} = \{a + bi : a, b \in \mathbb{R}\} \) is saying \( \mathbb{C} \) is the set whose elements are \( a + bi \), where \( a \) and \( b \) are real numbers. For example, in the notation above, \( 3 \in S, 2 \notin S, \) Jim’s football \( \in T \) but \( \notin U^+ \).

We haven’t really defined any of these sets rigorously. We have assumed that the reader understands what we mean. This suffices for the level of this course. A rigorous treatment is usually given in a logic course (usually for \( \mathbb{N} \)) or in analysis course for \( \mathbb{R} \). The set of real numbers can also be thought of as the set of all numbers written in a, possibly infinite, decimal expansion. Thus, \( 1, 2, 1/3 = 0.33333\ldots \) and indeed any rational number are elements of \( \mathbb{R} \) as are \( \pi = 3.1415926\ldots, e = 2.718281\ldots, \sqrt{2} = 1.414\ldots \) and so on.

We say that \( A \subseteq B \) (\( A \) is contained in \( B \)), or simply \( A \subset B \), if every element of \( A \) is an element of \( B \). For example \( \mathbb{N} \subset \mathbb{Z} \). We say that \( A = B \) if \( A \) is equal to \( B \), that is, if every element of \( A \) is an element of \( B \) and vice-versa. Note that \( A = B \) holds precisely when both \( A \subset B \) and \( B \subset A \).

The notation \( \emptyset \) stands for the empty set. It is a set but it has no elements. We let
\[ 
A \cap B = \{x : x \in A \text{ and } x \in B\} 
\]
be the intersection of \( A \) and \( B \), and
\[ 
A \cup B = \{x : x \in A \text{ or } x \in B\} 
\]
be the union of \( A \) and \( B \). For example, \( \{1, 3\} \cap \{n : n^2 = n\} = \{1\}, \mathbb{N} \cap \{x : -x \in \mathbb{N}\} = \{0\}, S_1 \cap T_3 = \emptyset \).

We shall also need arbitrary unions and intersections. Let \( I \) be a set (thought of as index set) and suppose for each \( i \in I \) we are given a set \( A_i \). Then
\[ 
\bigcap_{i \in I} A_i = \{x : x \in A_i, \forall i\}, 
\]
(\( \forall \) means “for all”) and
\[ 
\bigcup_{i \in I} A_i = \{x : x \in A_i, \text{ for some } i\}. 
\]

For example, define for \( i \in \mathbb{N} \),
\[ 
A_i = \{x \in \mathbb{N} : x \geq i\} 
\]
(so \( A_0 = \{0, 1, 2, 3, \ldots\}, A_1 = \{1, 2, 3, 4, \ldots\}, A_2 = \{2, 3, 4, 5, \ldots\} \) and so on). Then
\[ 
\bigcup_{i \in \mathbb{N}} A_i = \mathbb{N}, 
\]
while
\[ 
\bigcap_{i \in \mathbb{N}} A_i = \emptyset. 
\]
Here's a another example: for every real number \(x\), \(0 \leq x \leq 1\) define 
\[
S_x = \{(x, y) : 0 \leq y \leq x\}.
\]
Then 
\[
\bigcup_{0 \leq x \leq 1} S_x
\]
is the triangle in the plane whose vertices are \((0,0), (1,0), (1,1)\).

Yet another operation on sets is the \textit{difference} of sets: 
\[
A \setminus B = \{x : x \in A, x \notin B\}.
\]
We remark that \(A \cup B = B \cup A, A \cap B = B \cap A\) but usually \(A \setminus B\) is not equal to \(B \setminus A\).

Another definition we shall often use is that of the \textit{cartesian product}. Let \(A_1, A_2, \ldots, A_n\) be sets. Then 
\[
A_1 \times A_2 \times \cdots \times A_n = \{(x_1, x_2, \ldots, x_n) : x_i \in A_i\}.
\]
In particular, 
\[
A \times B = \{(a, b) : a \in A, b \in B\}.
\]

\textbf{Example 2.1.} Let \(A = \{1, 2, 3\}, B = \{1, 2\}\). Then 
\[
A \times B = \{(1,1), (1,2), (2,1), (2,2), (3,1), (3,2)\}.
\]
Note that \((3,1) \in A \times B\) but \((1,3) \notin A \times B\), because although \(1 \in A, 3 \notin B\).

\textbf{Example 2.2.} Let \(A = B = \mathbb{R}\). Then 
\[
A \times B = \{(x, y) : x \in \mathbb{R}, y \in \mathbb{R}\}.
\]
This is just the presentation of the plane in cartesian coordinates (which explains why we call such products “cartesian” products).

2.1. We now wish to make a general statements relating some of these operations. Once a statement is important enough to highlight it, it falls under the heading of Lemma, Proposition or Theorem (or, more colloquially, Claim, Assertion and so on). Usually, “Lemma” is reserved for technical statements often to be used in the proof of a proposition or a theorem. “Proposition” and “Theorem” are more or less the same. They are used for claims that are more conceptual, or central, with “Theorem” implying even more importance. However, none of these rules is absolute.

\textbf{Proposition 2.3.} Let \(I\) be a set. Let \(A\) be a set and \(B_i, i \in I\) be sets as well then 
\[
A \cap (\bigcup_{i \in I} B_i) = \bigcup_{i \in I} (A \cap B_i),
\]
and 
\[
A \cup (\bigcap_{i \in I} B_i) = \bigcap_{i \in I} (A \cup B_i).
\]
Furthermore, 
\[
A \setminus (\bigcup_{i \in I} B_i) = \bigcap_{i \in I} (A \setminus B_i),
\]
and

\[ A \setminus (\cap_{i \in I} B_i) = \cup_{i \in I}(A \setminus B_i). \]

3. Proofs: idea and technique

Proposition 2.3 is not obvious. It is not even clear at first sight whether it’s correct. For that reason we insist on proofs in mathematics. Proofs give us confidence that we are making true statements and they reveal, to a lesser or higher extent, why the statements hold true. The proof should demonstrate that the statements made are true. In fact, a few decades ago, people were using the word “demonstration” instead of “proof” and the famous QED in the end of the proof stands for “quod erat demonstrandum”, meaning, “which was to be demonstrated”. We shall prove some of the statements in the proposition now, leaving the rest as an exercise. Our method of proof for Proposition 2.3 is by a standard technique.

3.1. Proving equality by two inequalities. When one wants to show two real numbers \( x, y \) are equal. It is often easier to show instead that \( x \leq y \) and \( y \leq x \) and to conclude that \( x = y \).

In the same spirit, to show two sets \( A \) and \( B \) are equal, one may show that every element of \( A \) is an element of \( B \) and that every element of \( B \) is an element of \( A \). Thus, our principle of proof is

\[ A = B \]

if and only if

\[ x \in A \Rightarrow x \in B, \quad x \in B \Rightarrow x \in A. \]

(The notation \( \Rightarrow \) means “implies that”.)

Let us now prove the statement \( A \cap (\cup_{i \in I} B_i) = \cup_{i \in I}(A \cap B_i) \). The way you should write it in an assignment, test, or a research paper is:

**Proposition 3.1.** Let \( I \) be a set. Let \( A \) and \( B_i, i \in I, \) be sets then

\[ A \cap (\cup_{i \in I} B_i) = \cup_{i \in I}(A \cap B_i). \]

**Proof.** Let \( x \in A \cap (\cup_{i \in I} B_i) \) then \( x \in A \) and \( x \in \cup_{i \in I} B_i \). That is, \( x \in A \) and \( x \in B_{i_0} \) for some \( i_0 \in I \). Then \( x \in A \cap B_{i_0} \) and so \( x \in \cup_{i \in I}(A \cap B_i) \). We have shown so far that \( A \cap (\cup_{i \in I} B_i) \subset \cup_{i \in I}(A \cap B_i) \).

Conversely, let \( x \in \cup_{i \in I}(A \cap B_i) \). Then, there is some \( i_0 \in I \) such that \( x \in A \cap B_{i_0} \) and so for that \( i_0 \) we have \( x \in A \) and \( x \in B_{i_0} \). In particular, \( x \in A \) and \( x \in \cup_{i \in I} B_i \) and so \( x \in A \cap (\cup_{i \in I} B_i) \). □

(The □ designates that the proof is complete. It is equivalent to writing QED.)

Let us also do

\[ A \setminus (\cup_{i \in I} B_i) = \cap_{i \in I}(A \setminus B_i). \]
We use the same technique. Let \( x \in A \setminus (\cup_{i \in I} B_i) \) thus \( x \in A \) and \( x \notin \cup_{i \in I} B_i \). That means that \( x \in A \) and for all \( i \in I \) we have \( x \notin B_i \). That is, for all \( i \in I \) we have \( x \in A \setminus B_i \) and so \( x \in \cap_{i \in I} (A \setminus B_i) \).

Conversely, let \( x \in \cap_{i \in I} (A \setminus B_i) \). Then, for all \( i \in I \) we have \( x \in A \setminus B_i \). That is, \( x \in A \) and \( x \notin \cup_{i \in I} B_i \) and it follows that \( x \in A \setminus (\cup_{i \in I} B_i) \).

3.2. **Proof by contradiction and the contrapositive.** *Proof by contradiction* is a very useful technique, even though using it too often shows lack of deeper understanding of the subject. Suppose that some statement is to be proven true. In this technique one assumes that the statement is false and then proceeds to derive logical consequences of this assumption until an obvious contradiction arises. Here is an easy example that illustrates this:

**Claim.** There are no solution to the equation \( x^2 - y^2 = 1 \) in positive integers.

**Proof.** Assume not. Then there are positive integers \( x, y \) such that \( x^2 - y^2 = 1 \). Then, \((x - y)(x + y) = 1\). However, the only product of integers giving 1 is \(1 \times 1\) or \(-1 \times -1\) and, in any case, it follows that \( x - y = x + y \). It follows that \( 2y = (x + y) - (x - y) = 0 \) and so that \( y = 0 \). Contradiction. \(\square\)

Here is another example.

**Claim.** If \( x \) and \( y \) are two integers whose sum is odd, then exactly one of them is odd.

**Proof.** Suppose not. Then, either both \( x \) and \( y \) are odd, or both \( x \) and \( y \) are even. In the first case \( x = 2a+1, y = 2b+1 \), for some integers \( a, b \), and \( x+y = 2(a+b+1) \) is even – contradiction. In the second case, \( x = 2a, y = 2b \), for some integers \( a, b \), and \( x+y = 2(a+b) \) is even – contradiction again. \(\square\)

Somewhat related is the technique of *proving the contrapositive*. Let \( A \) and \( B \) be two assertions and let \( \neg A, \neg B \) be their negations. Logically the implication

\[ A \implies B \]

is equivalent

\[ \neg B \implies \neg A. \]

Here is an example. Let \( A \) be the statement “it rains” and \( B \) the statement “it’s wet outside”. Then \( \neg A \) is the statement “it doesn’t rain” and \( \neg B \) is the statement “it’s dry outside”. The meaning of \( A \implies B \) is “it rains therefore it’s wet outside” and its contrapositive is “it’s dry outside therefore it doesn’t rain”. Those statements are equivalent in the sense that each implies the other. Here is a mathematical example:

**Claim.** If \( x \) and \( y \) are integers such that \( xy \) is even then either \( x \) or \( y \) is even.
Proof. The contrapositive is: If both $x$ and $y$ are odd then $xy$ is odd. To prove that, write $x = 2a + 1, y = 2b + 1$ for some integers $a, b$. Then $xy = 4ab + 2a + 2b + 1 = 2(ab + a + b) + 1$, is one more than an even integer and so is odd. □

3.3. **Proof by Induction.** Induction is perhaps the most fun technique. Its logical foundations also lie deeper that the previous methods. The principle of induction, to be explained below, rests on the following **axiom**:

* **Axiom**: Every non-empty subset of $\mathbb{N}$ has a minimal element.

We remark that the axiom is actually intuitively obviously true. The reason we state it as an axiom is that when one develops the theory of sets in a very formal way from fundamental axioms the axiom stated above doesn’t follow from simpler axioms and, in one form or another, has to be included as an axiom.

**Theorem 3.2.** *(Principle of Induction)* Suppose that for every natural number $n \geq n_0$ we are given a statement $P_n$. Suppose that we know that:

1. $P_{n_0}$ is true.
2. If $P_n$ is true then $P_{n+1}$ is true.

Then $P_n$ is true for every $n$.

**Proof.** Suppose not. Then the set

$$S = \{ n \in \mathbb{N} : n \geq n_0, P_n \text{ is false} \}$$

is a non-empty set and therefore has a minimal element $a$. Note that $a > n_0$, because $P_{n_0}$ is true by (1), and so $a - 1 \geq n_0$. Now $a - 1 \not\in S$ because of the minimality of $a$ and so $P_{a-1}$ is true. But then, by (2), also $P_a$ is true. Contradiction. □

**Remark 3.3.** The mental picture I have is of a staircase. The first step is marked $n_0$ and each higher step by the integers $n_0 + 1, n_0 + 2, n_0 + 3, \ldots$. I know I can reach that first step and I know that if I can reach a certain step (say marked $n$) then I can reach the next one (marked $n + 1$). I conclude that I can reach every step.

**Example 3.4.** Prove that for every positive integer $n$ we have

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

(The statement is $P_n$ is $1 + 2 + \cdots + n = \frac{n(n + 1)}{2}$ and is made for $n \geq 1$, that is $n_0 = 1$.) The base case is when $n = 1$ (this is called the base case of the induction). In this case we need to show that $1 = \frac{1(1+1)}{2}$, which is obvious.

Now, we assume that statement true for $n$, that is we assume that

$$1 + 2 + \cdots + n = \frac{n(n + 1)}{2},$$
and we need to show it's true of \( n + 1 \). That is, we need to prove that

\[
1 + 2 + \cdots + n + (n + 1) = \frac{(n + 1)(n + 2)}{2}.
\]

(By achieving that we would have shown that if \( P_n \) is true then \( P_{n+1} \) is true.) We use the assumption that it’s true for \( n \) (that’s called the induction hypothesis) and write

\[
1 + 2 + \cdots + n + (n + 1) = \frac{n(n + 1)}{2} + n + 1
\]
\[
= \frac{n(n + 1) + 2(n + 1)}{2}
\]
\[
= \frac{n^2 + 3n + 2}{2}
\]
\[
= \frac{(n + 1)(n + 2)}{2}.
\]

Example 3.5. Here we prove the following statement: *Let \( q \neq 1 \) be a real number. Then for every \( n \in \mathbb{N} \) we have*

\[
1 + q + \cdots + q^n = \frac{1 - q^{n+1}}{1 - q}.
\]

The statement \( P_n \) is “for every real number \( q \neq 1 \), \( 1 + q + \cdots + q^n = \frac{1 - q^{n+1}}{1 - q} \).” The base case, that is the first \( n \) for the statement is being claimed true, is \( n = 0 \); the statement is then

\[
1 = \frac{1 - q}{1 - q},
\]

which is obviously true.

Now suppose that the statement is true for \( n \). That is,

\[
1 + q + \cdots + q^n = \frac{1 - q^{n+1}}{1 - q}.
\]

We need to show that

\[
1 + q + \cdots + q^{n+1} = \frac{1 - q^{n+2}}{1 - q}.
\]
Indeed,

\[ 1 + q + \cdots + q^{n+1} = (1 + q + \cdots + q^n) + q^{n+1} = \frac{1 - q^{n+1}}{1 - q} + q^{n+1} = \frac{1 - q^{n+1}}{1 - q} + \frac{(1 - q)q^{n+1}}{1 - q} = \frac{1 - q^{n+1}}{1 - q} + \frac{q^{n+1} - q^{n+2}}{1 - q} = \frac{1 - q^{n+2}}{1 - q}. \]

3.4. **Prove or disprove.** A common exercise (and a situation one often faces in research) is to prove or disprove a particular statement. For example,

“Prove or disprove: for every natural number \( n \), \( 4n + 1 \) is either a square or a sum of two squares.”

At that point you are requested first to form a hunch, a guess, an opinion about whether the statement is true or false. To form that hunch you can try some examples (\( 4 \times 0 + 1 = 1 = 1^2, 4 \times 1 + 1 = 5 = 1^2 + 2^2, 4 \times 2 + 1 = 9 = 3^2, 4 \times 3 + 1 = 13 = 2^2 + 3^2, \ldots \)) to see if the statement holds for these examples, see if it is similar to other statements you know to hold true, or, when at lost, throw the dice. After deciding on your initial position, if you believe the statement you should proceed to find a proof. If you don’t, then you have two options. You can try and show that if the statement is true it will imply a contradiction to a known fact, or to provide one counterexample. The statement being false doesn’t mean it’s false for every \( n \); it means it’s false of one \( n \). In the case in hand, if we take \( n = 5 \) we find that \( 4 \times 5 + 1 = 21 \) is neither a square nor a sum of squares (just try all possibilities) and so the statement is false.

3.5. **The pigeonhole principle.** The pigeonhole principle is deceptively simple, given how useful it is. It states the following:

*If there are more pigeons than pigeonholes then in one hole there must be at least two pigeons.*

**Example 3.6.** Let \( a, b, c, d, e, f \) be 6 integers. Then there are among them two integers whose difference is divisible by 5.

To prove this consider the residues of \( a, b, c, d, e, f \) upon division by 5. The residue is either 0, 1, 2, 3 or 4 (these are the 5 pigeonholes) but we get from our numbers 6 residues (those are the pigeons). Therefore, there are two numbers among \( a, b, c, d, e, f \) with the same residue. Say, \( a \) and \( b \) have the same residue, say \( r \). Then \( a - b \) is divisible by 5 (because \( a = 5a' + r, b = 5b' + r \) and so \( a - b = 5(a' - b') \)).
A similar example is the following:

**Example 3.7.** In any group of $n$ people, there are two that have the same number of friends within the group.

We prove that by induction on $n$. The cases $n = 1, 2$ are trivial. Suppose the claim for all $n \leq N$ and consider a group of $N+1$ people. If there is a person with 0 friends, we can look at the rest of the people. This is a group of $N$ people and the number of friends each has in this smaller group is the same as in the original group. We can apply induction to conclude that two have the same number of friends (initially in the smaller group, but in fact also in the larger group).

The other case is when each person in the group of $N+1$ people has at least 1 friend. By considering the number of friends each person in the group has, we get $N+1$ numbers between 1 and $N$ and so two must be equal.

4. Functions

There are more formal and less formal ways to define a function. Here we take the most pedestrian approach. Let $A$ and $B$ be sets. A function $f$ from $A$ to $B$,

$$f : A \rightarrow B,$$

is a rule assigning to each element of $A$ a unique element of $B$. The set $A$ is called the source, or the domain, of the function, and $B$ the target of the function. For $a \in A$, $f(a)$ is called the image of $a$ (under $f$) and $f(A) = \{f(a) : a \in A\}$ is the image of $f$.

**Example 4.1.** The simplest example is the identity function. Let $A$ be any set and define

$$1_A : A \rightarrow A$$

to be the function sending each element to itself. Namely,

$$1_A(x) = x,$$

for any $x \in A$.

**Example 4.2.** Let $A = \{1, 2, 3\}$, $B = \{1, 2\}$ and consider the following rules for $f : A \rightarrow B$.

1. $f(1) = 2$, $f(2) = 1$, $f(3) = 1$.
2. $f(1) = 1$ or 2, $f(2) = 2$, $f(3) = 1$.
3. $f(1) = 1$, $f(2) = 1$.

The first recipe defines a function from $A$ to $B$. The second recipe does not, because 1 is assigned two possible values. The third also doesn’t define a function because no information is given about $f(3)$.

**Example 4.3.** Consider

1. $f : \mathbb{R} \rightarrow \mathbb{R}$, $f(x) = \sqrt{x}$.
2. $f : \mathbb{R} \geq 0 \rightarrow \mathbb{R}$, $f(x) = y$, where $y$ is a real number such that $y^2 = x$. 
(3) $f : \mathbb{R}_{\geq 0} \to \mathbb{R}$, $f(x) =$ the non negative root of $x$.

The first definition fails because $-1$ doesn’t have a real root. The second definition fails because every positive number has 2 roots (differing by a sign) and it isn’t clear which root one is supposed to take. This problem also exists in the first definition. The third definition does define a function.

There are various ways to define a function. It could be done by writing down $f(a)$ for every $a \in A$ explicitly. It could be done by providing a formula and it could be done by giving some other description. For example, $A = B$ is the set of all people who ever lived, $f : A \to A$ is given by

$$f(a) = a’ s \text{ mother}.$$  

(This definitely looks like a good definition at first sight. However, the astute reader will note the problem here. If this function was truly well-defined then the set $A$ must be infinite, because if it were finite we would have a person who’s a descendant of itself. Since a mother is older from any of her childs by at least one year (to be safe) it follows, if $A$ and $f$ are well defined, that people have existed forever. The ways to resolve the paradox are rather amazing and I leave it to you as an amusing exercise.)

Here is some more notation: the symbol $\forall$ means “for all”. the symbol $\exists$ means “exists”. The symbol $\exists!$ means “exists unique”. A function can also be described as a set

$$\Gamma \subset A \times B,$$

with the following property: $\forall a \in A, \exists b \in B$ such that $(a, b) \in \Gamma$. We then define $f(a)$ to be the unique $b$ such that $(a, b) \in \Gamma$. Conversely, given a function $f$ we let

$$\Gamma = \Gamma_f = \{(a, f(a)) : a \in A\}.$$  

The set $\Gamma_f$ is called the graph of $f$.

**Example 4.4.** Let $A$ be a set and $\Gamma \subset A$ the “diagonal”,

$$\Gamma = \{(x, x) : x \in A\}.$$  

The function defined by $\Gamma$ is $1_A$.

**Example 4.5.** Let $f_1, f_2 : A \to B$ be functions. We say that $f_1 = f_2$ if for every $a \in A$ we have $f_1(a) = f_2(a)$. Equivalently (exercise!) if $\Gamma_{f_1} = \Gamma_{f_2}$.

4.1. **Injective, surjective, bijective, inverse image.** We introduce some attributes of functions. Let

$$f : A \to B$$

be a function. Then:

1. $f$ is called *injective* if $f(a) = f(a’) \Rightarrow a = a’$. (I.e., different elements of $A$ go to different elements of $B$.) Such a function is also called *one-one*.

2. $f$ is called *surjective* (or *onto*) if $\forall b \in B, \exists a \in A$ such that $f(a) = b$. (I.e., every element in the target is the image of some element in the source.)
(3) $f$ is called 	extit{bijective} if it is both injective and surjective. In that case, every element of $B$ is the image of a unique element of $A$.

Let $f: A \to B$ be a function. Let $U \subset B$. We define the pre-image of $U$ to be the set

$$f^{-1}(U) = \{ a : a \in A, f(a) \in U \}.$$ 

If $U$ consists of a single element, $U = \{u\}$, we also write $f^{-1}(u)$ and call it the fibre of $f$ over $u$.

**Example 4.6.**

1. $f: \mathbb{R} \to \mathbb{R}, f(x) = x^2$. Then $f$ is neither surjective (a square is always non-negative) nor injective ($f(x) = f(-x)$).

   We have $f^{-1}([1,4]) = [-2, -1] \cup [2, 4]$ and $f^{-1}(0) = \{0\}, f^{-1}(-1) = \emptyset$.

2. $f: \mathbb{R} \to \mathbb{R}_{\geq 0}, f(x) = x^2$. Then $f$ is surjective but not injective.

3. $f: \mathbb{R}_{\geq 0} \to \mathbb{R}_{\geq 0}, f(x) = x^2$. Then $f$ is bijective.

### 4.2. Composition of functions.

Let $f: A \to B, g: B \to C$, be functions. We define their composition, $g \circ f$, to be the function:

$$g \circ f : A \to C, \quad (g \circ f)(x) = g(f(x)).$$

**Lemma 4.7.** We have the following properties:

1. If $g \circ f$ is injective then $f$ is injective.
2. If $g \circ f$ is surjective then $g$ is surjective.

**Proof.** Suppose that $g \circ f$ is injective. Let $a, a' \in A$ be elements such that $f(a) = f(a')$. We need to show that $a = a'$. We have $g(f(a)) = g(f(a'))$ or otherwise said, $(g \circ f)(a) = (g \circ f)(a')$. Since $g \circ f$ is injective, $a = a'$.

Suppose now that $g \circ f$ is surjective. Let $c \in C$. We need to show that there is an element $b \in B$ with $g(b) = c$. Since $g \circ f$ is surjective, there is $a \in A$ such that $(g \circ f)(a) = c$. Let $b = f(a)$ then $g(b) = g(f(a)) = (g \circ f)(a) = c$. \[\square\]

### 4.3. The inverse function.

Let $f: A \to B$ be a bijective function. In this case we can define the inverse function

$$f^{-1}: B \to A,$$

by the property

$$f^{-1}(a) = b \text{ if } f(b) = a.$$ 

This is well defined: since $f$ is surjective such $b$ exists for every $a$ and is unique (because $f$ is injective). Thus $f^{-1}$ is a function. It is easy to verify that

$$f^{-1} \circ f = 1_A, \quad f \circ f^{-1} = 1_B.$$
5. Cardinality of a set

Imagine a group of students about to enter a lecture hall. The instructor wants to know if there are sufficient chairs for all the students. There are two ways to do that. One is to count both students and chairs separately and determine which is larger. The other is to ask each student to take a seat. If there are students left standing, the number of chairs is too small. If there are chairs left unoccupied, there are more chairs than students. In the remaining case there is a perfect match between students and chairs and so their number (cardinality) is equal.

This idea proves very powerful in discussing the cardinality (“size”, “magnitude”, “number of elements”) of sets, finite or infinite.

George Cantor has revolutionized mathematics, and human thought, by defining two sets $A, B$ possibly infinite to be of equal cardinality, denoted $|A| = |B|$, if there is a bijective function $f : A \rightarrow B$.

(Note that then there is an inverse function $f^{-1} : B \rightarrow A$ which is also a bijection, so it doesn’t matter if we require a bijection from $A$ to $B$ or from $B$ to $A$.) He defined the cardinality of $A$ to be no larger than $B$’s if there is an injective function $f : A \rightarrow B$.

We also say that the cardinality of $A$ is smaller than $B$’s if $|A| \leq |B|$ and $|A| \neq |B|$. As a sanity check we’d like to know that

**Proposition 5.1.** Let $A, B, C$ be sets. If $|A| = |B|$ and $|B| = |C|$ then $|A| = |C|$.

**Proof.** Let $f : A \rightarrow B, g : B \rightarrow C$ be bijections. Then

$$g \circ f : A \rightarrow C$$

is also a bijection. Indeed: if for some $x, y \in A$ we have $(g \circ f)(x) = (g \circ f)(y)$ then $g(f(x)) = g(f(y))$. Since $g$ is injective, $f(x) = f(y)$ and, since $f$ is injective, $x = y$.

To show $g \circ f$ is surjective, let $c \in C$ and choose $b \in B$ such that $g(b) = c$; such $b$ exists since $g$ is surjective. Since $f$ is surjective, there is an $a \in A$ such that $f(a) = b$. Then $(g \circ f)(a) = g(f(a)) = g(b) = c$. \[\square\]

To show this makes sense at all, we need to know the following

**Theorem 5.2** (Cantor-Bernstein). If $|A| \leq |B|$ and $|B| \leq |A|$ then $|A| = |B|$.

Although it does not require any sophisticated mathematics, the proof will not be given in this course. It is an ingenious and intricate proof. We would like to explain though why the theorem is not obvious.

What we are given that there is some injective function $f$ from $A$ to $B$ and some injective function $g$ from $B$ to $A$. We should conclude from that there is some bijective function $h$ from $A$ to $B$. It is not true that $h = f$ or $g$. One should somehow construct it from $f$ and $g$. Here is an
example: let $A$ be the set of points in the plane in distance at most 1 from the origin (the closed unit disk) and $B$ the square $[-1,1] \times [-1,1]$. The function

$$f : A \to B, \quad f(x) = x,$$

is a well-defined injective function but not a bijection. The function

$$g : B \to A, \quad g(b) = b/\sqrt{2},$$

is also a well-defined injective function, but not bijection. One can find a bijection from $A$ to $B$, but it is not immediately clear how to find it based on the knowledge of $f$ and $g$ (and in fact in this particular example, it is better to “rethink the situation” than to deduce it from $f$ and $g$).

A set $A$ is called countable (or enumerable) if it has the same cardinality as $\mathbb{N}$. In this case, there is a bijective function $f : \mathbb{N} \to A$. The elements of $A$ are thus $\{f(0), f(1), f(2), f(3), \ldots \}$. If we introduce the notation $a_i = f(i)$ then we can also enumerate the elements of $A$ as $\{a_0, a_1, a_2, a_3, \ldots \}$ and this explains the terminology.

**Example 5.3.** Let $A$ be the set $\{0,2,4,6,\ldots \}$ and $B$ the set $\{1,4,9,16,25,\ldots \}$. Then

$$|\mathbb{N}| = |A| = |B|.$$  

Indeed, one verifies that the functions

$$f : \mathbb{N} \to A, \quad f(x) = 2x,$$

and

$$g : \mathbb{N} \to B, \quad g(x) = x^2,$$

are bijections. We can then conclude that $|A| = |B|$ (though it is not hard either to write a bijection from $A$ to $B$).

**Example 5.4.** The cardinality of $\mathbb{N}$ is the cardinality of $A = \{x \in \mathbb{N} : x \text{ is not a square}\}$.

Instead of trying to write a bijection explicitly, which is not that straight-forward, we use the theorem. The function

$$f : A \to \mathbb{N}, \quad f(a) = a,$$

is injective. Thus, $|A| \leq |\mathbb{N}|$. Consider the function

$$f : \mathbb{N} \to A, \quad f(x) = 4x + 3.$$

First, this is well defined. Namely, $4x + 3$ is really in $A$. This is because if $n$ is square, $n$ leaves residue 1 or 0 when divided by 4 (if $n = (2m)^2 = 4m^2$ the residue is zero; if $n = (2m + 1)^2 = 4(m^2 + m) + 1$ the residue is one). But $4x + 3$ leaves residue 3. Clearly $f$ is injective and so $|\mathbb{N}| \leq |A|$.

**Proposition 5.5.** $|\mathbb{N}| = |\mathbb{Z}|$. 

Proof. We define
\[ f : \mathbb{Z} \rightarrow \mathbb{N}, \quad g(x) = \begin{cases} 2x & x \geq 0 \\ -2x - 1 & x < 0. \end{cases} \]
Then \( f \) is a bijective function, as is easy to check. \( \square \)

**Proposition 5.6.** \(|\mathbb{N}| = |\mathbb{N} \times \mathbb{N}|.\)

**Proof.** Define \( f : \mathbb{N} \rightarrow \mathbb{N} \times \mathbb{N} \) by \( f(n) = (n, 0) \). This is an injective function. Define \( g : \mathbb{N} \times \mathbb{N} \rightarrow \mathbb{N} \), \( g(n, m) = 2^n 3^m \).

This is also an injective function. If \( 2^n 3^m = 2^a 3^b \) then \( n = a, m = b \) by unique factorization (to be discussed in §10). We conclude that \(|\mathbb{N}| = |\mathbb{N} \times \mathbb{N}|.\) \( \square \)

**Corollary 5.7.** \(|\mathbb{Z}| = |\mathbb{Z} \times \mathbb{Z}|.\)

**Proof.** Let \( h : \mathbb{N} \rightarrow \mathbb{N} \times \mathbb{N} \) be a bijection. A bijection \( f : \mathbb{N} \rightarrow \mathbb{Z} \) induces a bijection \( g = (f, f) : \mathbb{N} \times \mathbb{N} \rightarrow \mathbb{Z} \times \mathbb{Z} \) and the composition
\[ \mathbb{Z} \xrightarrow{f} \mathbb{N} \xrightarrow{h} \mathbb{N} \times \mathbb{N} \xrightarrow{g^{-1}} \mathbb{Z} \times \mathbb{Z} , \]
is also a bijection. \( \square \)

**Exercise 5.8.** Prove that \(|\mathbb{N}| = |\mathbb{Q}|.\) (Hint: there’s an easy injection \( \mathbb{Q} \rightarrow \mathbb{Z} \times \mathbb{Z} \)).

When Cantor has laid down the foundations for the study of infinite cardinals he also dropped a bombshell:

**Theorem 5.9.** \(|\mathbb{N}| \neq |\mathbb{R}|.\)

The argument in the proof became known as Cantor’s diagonal argument and is used in many proofs.

**Proof.** Suppose that \(|\mathbb{N}| = |\mathbb{R}|.\) We can then enumerate the real numbers as \( a_0, a_1, a_2, \ldots \). Let us write the decimal expression of each number as
\[
a_0 = b_0^0 \ldots b_{n(0)}^0 \cdot c_0^0 c_1^0 c_2^0 c_3^0 \ldots \\
a_1 = b_1^1 \ldots b_{n(1)}^1 \cdot c_0^1 c_1^1 c_2^1 c_3^1 \ldots \\
a_2 = b_2^2 \ldots b_{n(2)}^2 \cdot c_0^2 c_1^2 c_2^2 c_3^2 \ldots \\
a_3 = b_3^3 \ldots b_{n(3)}^3 \cdot c_0^3 c_1^3 c_2^3 c_3^3 \ldots 
\]
where we agree to use \( 000000000 \ldots \) instead of \( 999999999 \ldots \) (so we write \( 1.0000000000 \ldots \) and not \( 0.9999999999 \ldots \), etc.).

Now consider the number
\[
0.e_0 e_1 e_2 e_3 e_4 \ldots , \quad e_i = \begin{cases} 3 & c_i^i \neq 3 \\ 4 & c_i^i = 3. \end{cases}
\]
This is a real number that differs from each $a_i$ at the $i + 1$-th digit after the decimal dot and hence is not equal to any $a_i$. It follows that the list $a_0, a_2, a_2, \ldots$ cannot consist of all real numbers and so we arrive at a contradiction. □

6. Number systems

Again we start with an apology of a sort. The formal discussion of number systems is a rather involved piece of mathematics. Our approach is pragmatic. We assume that at some level we all know what are integers and real numbers and no confusion shall arise there. We use those to define more complicated notions.

As we have already said, we denote the natural numbers by

$$\mathbb{N} = \{0, 1, 2, \ldots\}, \quad \mathbb{N}^+ = \{1, 2, 3, \ldots\}.$$  

We also denote the integers by

$$\mathbb{Z} = \{\ldots, -2, -1, 0, 1, 2, \ldots\}.$$  

The rational numbers are the set

$$\mathbb{Q} = \left\{ \frac{a}{b} : a, b \in \mathbb{Z}, b \neq 0 \right\}.$$  

The real numbers $\mathbb{R}$ are the “points on the line”. Each real number has a decimal expansion such as $0.19874526348 \ldots$ that may or may not repeat itself from some point on. For example:

$$\frac{1}{3} = 0.3333333 \ldots, \frac{1}{2} = 0.5000000 \ldots, \frac{1}{7} = 0.142857142857142857142857 \ldots$$  

It is a fact that a number is rational if and only if from some point on its decimal expansion becomes periodic.

The complex numbers are defined as the set

$$\mathbb{C} = \{a + bi : a, b \in \mathbb{R}\}.$$  

Here $i$ is a formal symbol. We can equally describe the complex numbers as points $(a, b) \in \mathbb{R}^2$ - the plane. The function $f : \mathbb{C} \to \mathbb{R}^2, f(a + bi) = (a, b)$ is bijective. The $x$-axis are now called the real axis and the $y$-axis the imaginary axis. If $z = a + bi$ is a complex number $a$ is called the real part of $z$ and is denoted $\text{Re}(z)$ and $b$ is called the imaginary part of $z$ and is denoted $\text{Im}(z)$. We therefore have

$$z = \text{Re}(z) + \text{Im}(z)i.$$  

The point corresponding to $z$ in the plane model is $(\text{Re}(z), \text{Im}(z))$.

One can perform arithmetic operations with the complex numbers using the following definitions:

$$-(a + bi) = -a - bi, \quad (a + bi) + (c + di) = (a + c) + (b + d)i.$$
Up to this point things look nice in the plane model as well:

\[-(a, b) = (-a, -b), \quad (a, b) + (c, d) = (a + c, b + d).\]

(This is the addition of the vectors.) The key point is that we can also define multiplication. The definition doesn’t have any prior interpretation in the model of the plane. We let

\[(a + bi)(c + di) = ac - bd + (ad + bc)i.\]

In particular,

\[i^2 = -1.\]

This shows that we have really gone beyond the realm of real numbers because there is no real number whose square is \(-1\) (or just negative, for that matter). The operations described above satisfy the usual rules of arithmetic, such as

\[(z + z') + z'' = z + (z' + z''), \quad z(z' + z'') = zz' + zz''.\]

(We shall later say that the complex numbers form a field.)

Let \(z = a + bi\) be a complex number. We define its complex conjugate of \(z\), \(\bar{z}\), as follows:

\[\bar{z} = a - bi.\]

**Lemma 6.1.** The complex conjugate has the following properties:

1. \(\bar{\bar{z}} = z\).
2. \(\bar{z_1 + z_2} = \bar{z}_1 + \bar{z}_2, \quad \bar{z_1 \cdot z_2} = \bar{z}_1 \cdot \bar{z}_2.\)
3. \(\text{Re}(z) = \frac{z + \bar{z}}{2}, \quad \text{Im}(z)i = \frac{z - \bar{z}}{2}\).
4. Define for \(z = a + bi\),

\[|z| = \sqrt{a^2 + b^2}.\]

(This is just the distance of the point \((a, b)\) from the origin.) Then \(|z|^2 = z \cdot \bar{z}\) and the following holds:

\[|z_1 + z_2| \leq |z_1| + |z_2|, \quad |z_1 \cdot z_2| = |z_1| \cdot |z_2|\).

**Proof.** Denote \(z = z_1 = a + bi, z_2 = c + di\).

We have \(\bar{z} = a - bi\) and so \(\bar{\bar{z}} = a + bi = z\). That is (1). For (2) we calculate

\[\bar{z_1 + z_2} = \frac{(a + c) + (b + d)i}{2} = \frac{(a + c) - (b + d)i}{2}
= \frac{a - bi + c - di}{2}
= \frac{a + bi + c + di}{2}
= \bar{z}_1 + \bar{z}_2.\]
Similarly,
\[
\frac{z_1z_2}{\sqrt{z_1z_2}} = \frac{(ac - bd) + (ad + bc)i}{\sqrt{(ac - bd) - (ad + bc)i}} = \frac{(a - bi)(c - di)}{a + bi \cdot c + di} = \frac{a^2 + b^2}{a^2 - b^2} = \frac{z_1 \cdot \bar{z}_2}{|z_1|^2 \cdot |z_2|^2}.
\]

We have \((z + \bar{z})/2 = ((a + bi) + (a - bi))/2 = a = \text{Re}(z)\) and \((z - \bar{z})/2 = ((a + bi) - (a - bi))/2 = bi = \text{Im}(z)i\), which is (3). Next, \(|z|^2 = a^2 + b^2 = (a + bi)(a - bi) = z \cdot \bar{z}\). Now,

\[
|z_1z_2|^2 = z_1z_2 \cdot \frac{\bar{z}_1\bar{z}_2}{z_1z_2} = z_1z_2 \cdot \frac{\bar{z}_1}{z_2} = z_1 \cdot \bar{z}_1 \cdot z_2 \cdot \bar{z}_2 = |z_1|^2 \cdot |z_2|^2.
\]

Thus, the assertion \(|z_1 \cdot z_2| = |z_1| \cdot |z_2|\) follows by taking roots. The inequality \(|z_1 + z_2| \leq |z_1| + |z_2|\) viewed in the plane model for complex numbers is precisely the assertion that the sum of the lengths of two sides of a triangle is greater or equal to the length of the third side. \(\square\)

**Example 6.2.** If \(z \neq 0\) then \(z\) has an inverse. Indeed, \(z \cdot \frac{\bar{z}}{|z|^2} = \frac{z \cdot \bar{z}}{|z|^2} = 1\). We write

\[
z^{-1} = \frac{\bar{z}}{|z|^2}.
\]

Just to illustrate we calculate the expression \(1 + 2i + \frac{3 - i}{13 + 5i}\). We have \(\frac{1}{13 + 5i} = \frac{1 - 5i}{26}\) and so \(\frac{3 - i}{13 + 5i} = (3 - i)(1 - 5i)/26 = -\frac{1}{13} - \frac{8}{13}i\) and thus \(1 + 2i + \frac{3 - i}{13 + 5i} = \frac{12}{13} + \frac{18}{13}i\).

6.1. **The polar representation.** Considering \(z = a + bi\) in the plane model as the vector \((a, b)\) we see that we can describe each complex number by the length \(r\) of the vector and the angle \(\theta\) it forms with the real axis. We have

\[
r = |z|, \quad \sin \theta = \frac{\text{Im}(z)}{|z|}, \quad \cos \theta = \frac{\text{Re}(z)}{|z|}.
\]

**Lemma 6.3.** If \(z_1\) has parameters \(r_1, \theta_1\) and \(z_2\) has parameters \(r_2, \theta_2\) then \(z_1z_2\) has parameters \(r_1r_2, \theta_1 + \theta_2\) (up to multiples of 360°).
Proof. We have \( r_1 r_2 = |z_1| |z_2| = |z_1 z_2| \) and this shows that \( r_1 r_2 \) is the length of \( z_1 z_2 \). Let \( \theta \) be the angle of \( z_1 z_2 \) then

\[
\sin \theta = \frac{\text{Im}(z_1 z_2)}{|z_1 z_2|} = \frac{\text{Re}(z_1) \text{Im}(z_2) + \text{Re}(z_2) \text{Im}(z_1)}{|z_1| |z_2|} = \frac{\text{Re}(z_1) \text{Im}(z_2)}{|z_1| |z_2|} + \frac{\text{Re}(z_2) \text{Im}(z_1)}{|z_1| |z_2|} = \cos \theta_1 \sin \theta_2 + \cos \theta_2 \sin \theta_1 = \sin(\theta_1 + \theta_2).
\]

Similarly, we get

\[
\cos \theta = \frac{\text{Re}(z_1 z_2)}{|z_1 z_2|} = \frac{\text{Re}(z_1) \text{Re}(z_2) - \text{Im}(z_1) \text{Im}(z_2)}{|z_1| |z_2|} = \cos(\theta_1) \cos(\theta_2) - \sin(\theta_1) \sin(\theta_2) = \cos(\theta_1 + \theta_2).
\]

It follows that \( \theta = \theta_1 + \theta_2 \) up to multiples of \( 360^0 \). \( \square \)

Let \( \theta \) be any real number. Let \( e^{i\theta} \) denote the unit vector whose angle is \( \theta \). That is, that complex number with length 1 and angle \( \theta \). Clearly we have

\[ e^{i\theta} = \cos \theta + i \sin \theta. \]

If \( z \) is any complex number with length \( r \) and angle \( \theta \) then we have the equality \( z = |z|e^{i\theta} \). The formula we have proven is

\[ z_1 z_2 = |z_1| |z_2| e^{i(\theta_1 + \theta_2)} \]

and in particular

\[ e^{i\theta_1} e^{i\theta_2} = e^{i(\theta_1 + \theta_2)}. \]

(Conversely, the identity \( e^{i\theta_1} e^{i\theta_2} = e^{i(\theta_1 + \theta_2)} \) implies the formula for multiplication). Let \( z = a + bi \) be a complex number then we define

\[ e^z = e^a e^{ib}, \]

where \( e^a \) is the usual exponential and \( e^{ib} \) is as defined above. We then have

\[ e^{z_1} e^{z_2} = e^{z_1 + z_2}. \]

We have defined here \( e^z \) in a purely formal way. One can show that for every complex number \( z \) the series

\[ 1 + z + \frac{z^2}{2!} + \frac{z^3}{3!} + \cdots + \frac{z^n}{n!} + \cdots \]
converges and is equal to $e^z$.

**Example 6.4.** Consider the polynomial $x^n - a = 0$, where $a$ is a non-zero complex number. We claim that this has $n$ distinct roots in $\mathbb{C}$. Write $a = re^{i\theta}$ (so $|a| = r$ and the line from 0 to $a$ forms an angle $\theta$ with the real axis). A complex number $z = Re^{i\Theta}$ is a solution to the equation if and only if $z^n = R^n e^{in\Theta} = re^{i\theta}$. That is, if and only if

$$R^n = r, \quad n\Theta \equiv \theta \pmod{2\pi}.$$  

Thus, the solutions are exactly

$$z = r^{1/n} e^{i(\theta + 2\pi j/n)}, \quad j = 0, 1, \ldots, n-1.$$  

In particular, taking $a = 1$ the solutions are called the *roots of unity* of order $n$. There are precisely $n$ of them. Those are the points on the unit circle having angle $0, \frac{1}{n} \cdot 2\pi, \frac{2}{n} \cdot 2\pi, \ldots, \frac{n-1}{n} \cdot 2\pi$.

**6.2. The Fundamental Theorem of Algebra.** A *complex polynomial* $f(x)$ is an expression of the form $a_n x^n + a_{n-1} x^{n-1} + \cdots + a_1 x + a_0$, where $n$ is a non-negative integer, $x$ is a variable, and the coefficients $a_i$ are complex numbers. (If all the coefficients are real we may call it a *real polynomial*; if all the coefficients are rational numbers we may call it a *rational polynomial* and so on. But note that $x^2 + 1$ is both a rational, a real and a complex polynomial.) The *zero polynomial*, denote 0, is the case when $n = 0$ and $a_0 = 0$.

A polynomial defines a function

$$f : \mathbb{C} \to \mathbb{C}, \quad z \mapsto f(z) = a_n z^n + a_{n-1} z^{n-1} + \cdots + a_1 z + a_0.$$  

(The notation $\mapsto$ means “maps to”.) If $a_n \neq 0$ then we say $f$ has *degree* $n$. If $f(z) = 0$ we say that the complex number $z$ is a *root* (or a *solution*, or a *zero*) of the polynomial $f$.

**Example 6.5.** Consider the polynomial $f(x) = x^2 + 1$. It has degree 2 and $f(i) = i^2 + 1 = -1 + 1 = 0$, $f(-i) = (-i)^2 + 1 = -1 + 1 = 0$. So $i$ and $-i$ are roots of $f$. (This is a special case of Example 6.4.)

**Theorem 6.6.** *(The Fundamental Theorem of Algebra)* Let $f(x)$ be a complex polynomial of degree at least 1. Then $f(x)$ has a root in $\mathbb{C}$.

The proof of the theorem is beyond the scope of this course. It has many proofs. In Honours Algebra 4 MATH371 one sees an algebraic proof; in Complex Variables and Transforms MATH 381 one sees an analytic proof.

**Proposition 6.7.** Let $f(x) = a_n x^n + a_{n-1} x^{n-1} + \cdots + a_1 x + a_0$ be a complex non-zero polynomial of degree $n$. Then

$$f(x) = a_n \prod_{i=1}^{n} (x - z_i),$$

where $z_1, \ldots, z_n$ are the roots of $f$. (This is a special case of Theorem 6.6.)
for suitable complex numbers, not necessarily distinct, \( z_i \). The numbers \( z_i \) are all roots of \( f \) and any root of \( f \) is equal to some \( z_i \). Moreover, this factorization is unique.

**Proof.** We prove the result by induction on \( n \). For \( n = 0 \) we understand the product \( \prod_{i=1}^{n}(x - z_i) \) as one (this is a convention: the empty product is equal to one, the empty sum is equal to zero.) and so the claim is just that a constant polynomial is equal to its leading coefficient. Clear.

Now, assume that \( f \) has degree at least one. By the Fundamental Theorem of Algebra there is a complex number \( z_n \) say such that \( f(z_n) = 0 \). We claim that for every complex number \( z \) we can write

\[
f(x) = (x - z)g(x) + r,
\]

where \( g(x) \) is a polynomial of degree \( n - 1 \) and leading coefficient \( a_n \) and \( r \) is a complex number. Indeed, write \( g(x) = b_{n-1}x^{n-1} + \cdots + b_1x + b_0 \) and equate coefficients in \( (x - z)g(x) = b_{n-1}x^n + (b_{n-2} - zb_{n-1})x^{n-1} + \cdots + (b_0 - zb_1)x \) and \( f(x) \). We want complex numbers \( b_0, \ldots, b_{n-1} \) such that

\[
b_{n-1} = a_n, (b_{n-2} - zb_{n-1}) = a_{n-1}, \ldots, (b_0 - zb_1) = a_1,
\]

and there is no problem solving these equations. Thus, we can choose \( g(x) \) with a leading coefficient \( a_n \) such that \( f(x) - (x - z)g(x) = r \) is a constant.

Now, apply that for \( z = z_n \). We have \( f(x) - (x - z_n)g(x) = r \). We view \( r \) as a polynomial and substitute \( x = z_n \). We get

\[
f(z_n) - (z_n - z_n)g(z_n) = r.
\]

Since \( f(z_n) = 0 \) we conclude that \( r = 0 \).

We showed that if \( f(z_n) = 0 \) then

\[
f(z_n) = (x - z_n)g(x), \quad g(x) = b_{n-1}x^{n-1} + \cdots + b_0.
\]

In fact, \( b_{n-1} = a_n \). Using the induction hypothesis, we have

\[
g(x) = a_n \prod_{i=1}^{n-1}(x - z_i),
\]

for some complex numbers \( z_i \) and so

\[
f(x) = a_n \prod_{i=1}^{n}(x - z_i).
\]

We note that \( f(z_j) = a_n \prod_{i=1}^{n}(z_j - z_i) = 0 \), because the product contains the term \((z_j - z_j)\). If \( f(z) = 0 \) then \( a_n \prod_{i=1}^{n}(z - z_i) = 0 \). But, if a product of complex numbers is zero one of the number is already zero. Since \( a_n \neq 0 \), we must have \( z = z_i \) for some \( i \).

It remains to prove the uniqueness of the factorization. Suppose that

\[
f(x) = a_n \prod_{i=1}^{n}(x - z_i) = a \prod_{i=1}^{n}(x - t_i).
\]
Since the leading coefficient of \( f \) is \( a_n \) we must have \( a = a_n \). We now argue by induction. The case of degree 0 is clear. Assume \( f \) has degree greater than zero. Then the \( t_i \) are roots of \( f \) and so \( t_1 \) is equal to some \( z_i \). But we may re-index the \( z_i \) so that \( t_1 = z_1 \). Dividing both sides by \( x - z_1 \) we then conclude that\(^2\)

\[
 a_n \prod_{i=2}^n (x - z_i) = a_n \prod_{i=2}^n (x - t_i),
\]

and, by induction, \( z_i = t_i \) for all \( i \). \(\square\)

We remark that for \( n = 1, 2 \) the result is well known:

\[
 ax + b = a \left( x - \frac{-b}{a} \right),
\]

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

There are also formulas for the roots for polynomials of degree 3 and 4, but in degrees 5 and higher no such formulas exist. This follows from Galois theory, taught in MATH371.

7. Fields and rings - definitions and first examples

In the examples we have already discussed, or the definitions we have made, there are implicit structures that we want to define now in a formal way. At this point we just provide the definitions and reconsider previous examples. Later we shall enter a systematic development of the theory.

An operation (more pedantically called a “binary operation”) on a set \( R \) is a function

\[
w : R \times R \to R.
\]

That is, it is a rule taking two elements of \( R \) and returning a new one. For example:

\[
w : \mathbb{C} \times \mathbb{C} \to \mathbb{C}, \quad w(z_1, z_2) = z_1 + z_2,
\]

or

\[
w : \mathbb{C} \times \mathbb{C} \to \mathbb{C}, \quad w(z_1, z_2) = z_1 z_2.
\]

Often, for a general set \( R \) we may denote \( w(z_1, z_2) \) by \( z_1 + z_2 \), or \( z_1 z_2 \), if we want to stress the fact that the operation behaves like addition, or multiplication.

**Definition 7.1.** A ring \( R \) is a non-empty set together with two operations, called “addition” and “multiplication” that are denoted, respectively, by

\[
(x, y) \mapsto x + y, \quad (x, y) \mapsto xy.
\]

One requires the following axioms to hold:

\(^2\)We say that \( f(x)/g(x) = h(x) \) if \( h(x) \) is a polynomial such that \( f(x) = g(x)h(x) \). We shall see later that \( h(x) \) is uniquely determined. In our case clearly \( f(x)/(x - z_1) = a_n \prod_{i=2}^n (x - z_i) \).
(1) \( x + y = y + x, \forall x, y \in R \). (Commutativity of addition)
(2) \( (x + y) + z = x + (y + z), \forall x, y, z \in R \). (Associativity of addition)
(3) There exists an element in \( R \), denoted 0, such that \( 0 + x = x, \forall x \in R \). (Neutral element for addition)
(4) \( \forall x \in R, \exists y \in R \) such that \( x + y = 0 \). (Inverse with respect to addition)
(5) \( (xy)z = x(yz), \forall x, y, z \in R \). (Associativity of multiplication)
(6) There exists an element \( 1 \in R \) such that \( 1 \cdot x = x, \forall x \in R \). (Neutral element for multiplication)
(7) \( z(x + y) = zx + zy, (x + y)z = xz + yz, \forall x, y, z \in R \). (Distributivity)

We remark that for us a ring always has an identity element with respect to multiplication. Not all authors insist on that, though.

Note that the multiplication is not assumed to be commutative.

**Definition 7.2.** If \( xy = yx \) for all \( x, y \in R \), we say \( R \) is a commutative ring. If for every non-zero \( x \in R \) there is an element \( y \in R \) such that \( xy = yx = 1 \), and also \( 0 \neq 1 \) in \( R \), we call \( R \) a division ring. A commutative division ring is called a field.

**Example 7.3.** \( \mathbb{Z} \) is a commutative ring. It is not a division ring and so it is not a field.

**Example 7.4.** The rational numbers \( \mathbb{Q} \) form a field. The real numbers \( \mathbb{R} \) form a field. In both those cases we are forced to assume the properties as “well known”. The complex numbers also form a field, in fact we have at some level already used all the axioms implicitly in our calculations, but now we prove it formally using that \( \mathbb{R} \) is a field.

**Proposition 7.5.** \( \mathbb{C} \) is a field.

**Proof.** Let \( z_1 = a_1 + b_1i, z_2 = a_2 + b_2i, z_3 = a_3 + b_3i \). We verify the axioms:
1. \( z_1 + z_2 = (a_1 + a_2) + (b_1 + b_2)i = (a_2 + a_1) + (b_2 + b_1)i = z_2 + z_1 \).
2. \( (z_1 + z_2) + z_3 = [(a_1 + a_2) + (b_1 + b_2)i] + a_3 + b_3i = [(a_1 + a_2) + a_3] + [(b_1 + b_2) + b_3i] = [a_1 + (a_2 + a_3)] + [b_1 + (b_2 + b_3)]i = z_1 + [(a_2 + a_3) + (b_2 + b_3)i] = z_1 + (z_2 + z_3) \).
3. Clearly 0 + \( z_1 = z_1 \).
4. We have \( (-a_1 - b_1)i + (a_1 + b_1i) = (-a_1 + a_1) + (-b_1 + b_1)i = 0 + 0i = 0 \).
5. \( (z_1z_2)z_3 = [(a_1 + b_1i)(a_2 + b_2i)](a_3 + b_3i) = ((a_1a_2 - b_1b_2) + (a_1b_2 + b_1a_2)i)(a_3 + b_3i) = (a_1a_2 - b_1b_2)a_3 - (a_1b_2 + b_1a_2)b_3 + ((a_1a_2 - b_1b_2)b_3 + (a_1b_2 + b_1a_2)a_3)i = a_1a_2a_3 - b_1b_2a_3 - a_1b_2a_3 + a_1a_2a_3i \). One now develops the product \( z_1(z_2z_3) \) in the same way and checks that the answers match. We don’t do that here.
6. Clearly \( 1 \cdot z_1 = z_1 \cdot 1 = z_1 \).
7. \( z_1(z_2 + z_3) = (a_1 + b_1i)((a_2 + a_3) + (b_2 + b_3)i) = a_1(a_2 + a_3) - b_1(b_2 + b_3) + (b_1(a_2 + a_3) + a_1(b_2 + b_3)i) = (a_1a_2 - b_1b_2) + (b_1a_2 + a_1b_2)i + (a_1a_3 - b_1b_3) + (b_1a_3 + a_1b_3)i = (a_1 + b_1i)(a_2 + b_2i) + (a_1 + b_1i)(a_3 + b_3i) = z_1z_2 + z_1z_3 \).

Before proving the next property, we check that \( z_1z_2 = z_2z_1 \). We have \( z_1z_2 = (a_1 + b_1i)(a_2 + b_2i) = a_1a_2 - b_1b_2 + (a_1b_2 + b_1a_2)i = a_2a_1 - b_2b_1 + (a_2b_1 + b_2a_1)i = (a_2 + b_2i)(a_1 + b_1i) = z_2z_1 \).

In particular, \( (z_1 + z_2)z_3 = z_3(z_1 + z_2) = z_3z_1 + z_3z_2 = z_1z_3 + z_2z_3 \).
Finally, as we have already seen, if \( z_1 \neq 0 \) then \( z_1 \cdot \frac{\bar{z}_1}{|z_1|^2} = 1 \). We proved that \( \mathbb{C} \) is a field. \( \square \)

7.1. **Some formal consequences of the axioms.** We note some useful formal consequences of the axioms defining a ring:

1. The element 0 appearing in axiom (3) is unique. Indeed, if \( q \) is another element with the same property then \( q + x = x \) for any \( x \) and in particular \( q + 0 = 0 \). But also, using the property of 0 and commutativity, we have \( q + 0 = 0 + q = q \). So \( q = 0 \).

2. The element \( y \) appearing in axiom (4) is unique. Indeed, if for a given \( x \) we have \( x + y = x + y' = 0 \) then \( y = y + (x + y') = (y + x) + y' = (x + y) + y' = 0 + y' = y' \). We shall denote \( y \) by \(-x\).

3. We have \(-(-x) = x \) and \(-(x + y) = -x - y\), where, technically \(-x - y\) means \(-x + (-y)\).

   To prove that, it is enough, after what we have just proven, to show that \(-x + x = 0\) and that \((x + y) + (-x - y) = 0\). But, \(-x + x = x + (-x) = 0\) and \((x + y) + (-x - y) = x + (-x) + y + (-y) = 0 + 0 = 0\).

4. The element 1 in axiom (6) is unique. (Use the same argument as in (1)).

5. We have \( x \cdot 0 = 0, 0 \cdot x = 0 \). Indeed, \( x \cdot 0 = x \cdot (0 + 0) = x \cdot 0 + x \cdot 0 \). Let \( y = x \cdot 0 \) then \( y = y + y \) and so \( 0 = -y + y = -y + (y + y) = (-y + y) + y = 0 + y = y \).

We shall see many examples of rings and fields in the course. For now, we just give one more definition and some examples.

**Definition 7.6.** Let \( R \) be a ring. A subset \( S \subset R \) is called a subring if \( 0, 1 \in S \) and if \( a, b \in S \) implies that \( a + b, -a, ab \in S \).

Note that the definition says that the operations of addition and multiplication in \( R \) give operations of addition and multiplication in \( S \) (namely, the outcome is in \( S \) and so we get functions \( S \times S \rightarrow S \)), satisfying all the axioms of a ring. It follows that \( S \) is a ring whose zero element is that of \( R \) and whose identity is that of \( R \).

For example, \( \mathbb{Q} \) is a subring of \( \mathbb{R} \), \( \mathbb{Z} \) is a subring of \( \mathbb{Q} \) and so on.

**Example 7.7.** Consider the set \( \{0, 1\} \) with the following addition and multiplication tables.

<table>
<thead>
<tr>
<th></th>
<th>0</th>
<th>1</th>
</tr>
</thead>
<tbody>
<tr>
<td>+</td>
<td>0</td>
<td>1</td>
</tr>
<tr>
<td>0</td>
<td>0</td>
<td>1</td>
</tr>
<tr>
<td>1</td>
<td>1</td>
<td>0</td>
</tr>
<tr>
<td>×</td>
<td>0</td>
<td>1</td>
</tr>
<tr>
<td>0</td>
<td>0</td>
<td>0</td>
</tr>
<tr>
<td>1</td>
<td>0</td>
<td>1</td>
</tr>
</tbody>
</table>

One can verify that this is a ring by directly checking the axioms. (We shall later see that this is the ring of integers modulo 2).

**Example 7.8.** Consider all expressions of the form \( \{a + b\sqrt{2} : a, b \in \mathbb{Z}\} \). We use the notation \( \mathbb{Z}[\sqrt{2}] \) for this set. This set is actually a ring. Since \( \mathbb{Z}[\sqrt{2}] \subset \mathbb{R} \) and \( \mathbb{R} \) is a ring (even a field!), it is enough to check it’s a subring. Indeed, \( 0, 1 \in \mathbb{Z}[\sqrt{2}] \). Suppose \( a + b\sqrt{2}, c + d\sqrt{2} \in \mathbb{Z}[\sqrt{2}] \).

Then: (1) \((a+b\sqrt{2})+(c+d\sqrt{2}) = (a+c)+(b+d)\sqrt{2} \in \mathbb{Z}[\sqrt{2}]\); (2) \(-(a+b\sqrt{2}) = -a-b\sqrt{2} \in \mathbb{Z}[\sqrt{2}]\);
(3) \((a+b\sqrt{2})(c+d\sqrt{2}) = (ac + 2bd) + (ad + bc)\sqrt{2} \in \mathbb{Z}[\sqrt{2}]\).
Let now
\[ \mathbb{Q}[\sqrt{2}] = \{ a + b\sqrt{2} : a, b \in \mathbb{Q} \}. \]
This is a field. The verification that this is a subring of \( \mathbb{C} \) is the same as above. It is thus a commutative ring in which \( 0 \neq 1 \). We need to show inverse for multiplication. If \( a + b\sqrt{2} \) is not zero then either \( a \) or \( b \) are not zero. If
\[ c = a^2 - 2b^2 \]
is zero then either \( b = 0 \) (but then \( a \neq 0 \) and so \( c \neq 0 \), so this case doesn’t happen), or \( \sqrt{2} = a/b \) is a rational number. We shall prove in Proposition 10.9 that this is not the case. Thus, \( c \neq 0 \). Now, \( \frac{a}{c} - \frac{b}{c} \sqrt{2} \in \mathbb{Q}[\sqrt{2}] \) and it is easy to check that
\[ (a + b\sqrt{2}) \left( \frac{a}{c} - \frac{b}{c} \sqrt{2} \right) = 1. \]
Part 2. Arithmetic in \( \mathbb{Z} \)

In this part of the course we are going to study arithmetic in the ring of integers \( \mathbb{Z} \). We are going to focus on particular properties. Our choice of properties is motivated by an analogy to be drawn later between integers and polynomials. In fact, there even more general classes of rings for which one can extend this analogy, called Euclidean rings, but in this course we shall not discuss them beyond providing some margin notes.

8. Division

Theorem 8.1. (Division with residue)\footnote{A commutative ring is called an integral domain if \( xy = 0 \) implies \( x = 0 \) or \( y = 0 \). An integral domain \( R \) is called a Euclidean ring if there is a function \( | \cdot | : R - \{0\} \rightarrow \mathbb{N} \) such that for all \( x, y \) in \( R \) with \( y \neq 0 \) there are elements \( q, r \in R \) such that \( x = qy + r \) and either \( r = 0 \) or \( |r| < |x| \). (But in general \( q, r \) are not unique.) Thus, the theorem we are proving is that \( \mathbb{Z} \) is a Euclidean ring.} Let \( a, b \) be integers with \( b \neq 0 \). There exist integers \( q, r \) such that

\[
a = qb + r, \quad 0 \leq r < |b|.
\]

Moreover, \( q \) and \( r \) are uniquely determined.

Proof. For simplicity, assume \( b > 0 \). Very similar arguments prove the case \( b < 0 \).

Consider the set

\[
S = \{a - bx : x \in \mathbb{Z}, a - bx \geq 0\}.
\]

\( S \) is a non-empty set, indeed, if \( a > 0 \) take \( x = 0 \) and \( a \in S \). If \( a < 0 \) take \( x = a \) and \( a - bx = a(1 - b) \geq 0 \) (because \( b > 0 \) and so \( b \geq 1 \)). That is, \( a(1 - b) \in S \). It follows that \( S \) has a minimal element \( r = a - bq \) for some \( q \). Then \( r < b \); otherwise, \( 0 \leq r - b = a - b(q + 1) \) is an element of \( S \) as well and smaller that \( r \), which is a contradiction. It follows that

\[
a = bq + r, \quad 0 \leq r < b.
\]

We now show that \( q \) and \( r \) are unique. Suppose

\[
a = bq' + r', \quad 0 \leq r' < b.
\]

If \( q = q' \) then also \( r = a - bq = a - bq' = r' \). Else, either \( q > q' \) or \( q' > q \). We then get

\[
0 = bq + r - (bq' + r') = b(q - q') + (r - r').
\]

If \( q > q' \) then \( r' = r + b(q - q') \geq r + b \geq b \). Contradiction. If \( q < q' \) we get \( r = r' + b(q' - q) \geq b \) and again a contradiction. \( \square \)
9. GCD and the Euclidean algorithm

**Definition 9.1.** Let $a, b$ be integers. We say that $a|b$ (read, $a$ divides $b$) if there is an element $c \in \mathbb{Z}$ such that $b = ac$.

Here are some properties:
1. $a|b \Rightarrow a|-b$.
2. $a|b \Rightarrow a|bd$ for any $d \in \mathbb{Z}$.
3. $a|b, a|d \Rightarrow a|(b + d)$.

**Proof.** Write $b = ac$. Then $-b = a \cdot (-c)$ and so $a|-b$. Also, $bd = a \cdot (cd)$ and so $a|bd$.

Write also $d = ae$. Then $b + d = a \cdot (c + e)$ and so $a|(b + d)$.

**Corollary 9.2.** Let $a \neq 0$. $a|b$ if and only if in dividing $b$ in $a$ with residue, $b = aq + r$, the residue $r$ is zero.

**Proof.** If the residue $r = 0$ then $b = aq$ and so $a|b$. If $a|b$ and $b = aq + r$ then $a|(b - aq)$, i.e., $a|r$. But $r < |a|$ and so that’s possible only if $r = 0$.

9.1. GCD.

**Definition 9.3.** Let $a, b$ be integers, not both zero. The greatest common divisor (gcd) of $a$ and $b$, denoted gcd$(a, b)$ or just $(a, b)$ if the context is clear, is the largest integer dividing both $a$ and $b$.

**Theorem 9.4.** Let $a, b$ be integers, not both zero, and $d = (a, b)$ their gcd. Then every common divisor of $a$ and $b$ divides $d$. There are integers $u, v$ such that

$$d = ua + vb.$$  

Moreover, $d$ is the minimal positive number that has the form $ua + vb$.

**Proof.** Let

$$S = \{ma + nb : m, n \in \mathbb{Z}, ma + nb > 0\}.$$  

First note that $S \neq \emptyset$. Indeed, $a + bb \in S$. Let $D$ be the minimal element of $S$. Then, for some $u, v \in \mathbb{Z}$ we have $D = ua + vb$.

We claim that $D = d$. To show $D|a$, write $a = qD + r, 0 \leq r < D$. Then, $D > r = a - qD = a - q(ua + vb) = (1 - qu)a - qvb$. If $r \neq 0$ then $r = (1 - qu)a - qvb$ is an element of $S$ smaller than $D$ and that’s a contradiction. It follows that $r = 0$, that is $D|a$. In the same way, $D|b$.

On the other hand, let $e$ be any common divisor of $a$ and $b$. Then $e$ also divides $ua + vb = D$. It follows that $D$ is the largest common divisor of $a, b$, so $D = d$, and also that any other common divisor divides it.

**Corollary 9.5.** If $a|bc$ and $gcd(a, b) = 1$ then $a|c$.

**Proof.** We have $1 = ua + vb$ for some integers $u, v$. Since $a|uac$ and $a|vbc$ we have $a|uac + vbc = c$.  

\[\Box\]
9.2. **The Euclidean algorithm.** The question arises: how do we compute in practice the gcd of two integers? This is a very practical issue, even in the simple task of simplifying fractions! As we shall see, there are two methods. One method uses the prime factorization of the two numbers – we shall discuss that later. The other method, *which is much more efficient*, is the Euclidean algorithm.

**Theorem 9.6.** *(The Euclidean Algorithm)* Let $a, b$ be positive integers with $a \geq b$. If $b|a$ then $\gcd(a, b) = b$. Else perform the following recursive division with residue:

\[
\begin{align*}
    a &= bq_0 + r_0, \quad 0 < r_0 < b, \\
    b &= r_0q_1 + r_1, \quad 0 \leq r_1 < r_0 \\
    r_0 &= r_1q_2 + r_2, \quad 0 \leq r_1 < r_2 \\
    &\vdots
\end{align*}
\]

For some $t$ we must first get that $r_{t+1} = 0$. That is,

\[
\begin{align*}
    r_{t-2} &= r_{t-1}q_t + r_t, \quad 0 \leq r_t < r_{t-1} \\
    r_{t-1} &= r_{tq_{t+1}}.
\end{align*}
\]

Then $r_t$ is the gcd of $a$ and $b$.

Before giving the proof we provide two examples.

1). Take $a = 113, b = 54$. Then

\[
\begin{align*}
    113 &= 54 \cdot 2 + 5 \\
    54 &= 5 \cdot 10 + 4 \\
    5 &= 4 \cdot 1 + 1 \\
    4 &= 4 \cdot 1.
\end{align*}
\]

Thus $\gcd(113, 54) = 1$.

2). Now take $a = 442, b = 182$. Then

\[
\begin{align*}
    442 &= 182 \cdot 2 + 78 \\
    182 &= 78 \cdot 2 + 26 \\
    78 &= 26 \cdot 3
\end{align*}
\]

and so $\gcd(442, 182) = 26$.

**Proof.** Let $d = \gcd(a, b)$. We claim that $d|r_n$ for every $n$. We prove that by induction: First, $d|a, d|b$ then $d|(a - bq_0) = r_0$. Suppose that $d|r_i, i = 0, 1, 2, \ldots, n$. Since $r_{n+1} = r_{n-1} - r_nq_{n+1}$ we get that $d|r_{n+1}$ as well. In particular, $d|r_t$.

We now show that $r_t|a, r_t|b$. It then follows that $r_t|d$ and therefore $r_t = d$. We again prove that by induction. We have $r_t|r_t$ and $r_t|r_tq_{t+1} = r_t - 1$. Suppose we have already shown that $r_t$ divides $r_t, r_{t-1}, \ldots, r_n$. Then, since $r_{n-1} = r_{n-1} - r_nq_{n+1}$ we also get $r_t|r_{n-1}$. Therefore, $r_t$ divides $r_0, r_1, \ldots, r_t$. Again, $b = r_0q_1 + r_1$ and so $r_t|b$ and then $a = bq_0 + r_0$ and so $r_t|a$.

A further bonus supplied by the Euclidean algorithm is that it allows us to find $u, v$ such that $\gcd(a, b) = ua + vb$. We just illustrate it in two examples:

1). Take $a = 113, b = 54$. Then, as we saw,

\[
\begin{align*}
    113 &= 54 \cdot 2 + 5 \\
    54 &= 5 \cdot 10 + 4
\end{align*}
\]
5 = 4 \cdot 1 + 1
4 = 4 \cdot 1
and so \( \gcd(113, 54) = 1 \). We have 1 = 5 − 4 \cdot 1, and we keep substituting for the residues we now have, expressions using previous residues (the important numbers to modify are the residues not the quotients \( q_i \)). 4 = 54 − 5 \cdot 10 and we get 1 = 5 − (54 − 5 \cdot 10) = −54 + 5 \cdot 11. Next, 5 = 113 − 54 \cdot 2 and we get 1 = −54 + 5 \cdot 11 = −54 + (113 − 54 \cdot 2) \cdot 11 = 54 \cdot (−23) + 113 \cdot 11. Thus,

\[ 1 = \gcd(54, 113) = −23 \cdot 54 + 11 \cdot 113. \]

2). Now take \( a = 442, b = 182 \). Then

\[
\begin{align*}
442 &= 182 \cdot 2 + 78 \\
182 &= 78 \cdot 2 + 26 \\
78 &= 26 \cdot 3 \\
\end{align*}
\]
and so \( \gcd(442, 182) = 26 \). Here the process is easier: 26 = 182 − 78 \cdot 2 = 182 − (442 − 182 \cdot 2) \cdot 2 = 5 \cdot 182 − 2 \cdot 442.

\[ 26 = \gcd(182, 442) = 5 \cdot 182 − 2 \cdot 442. \]

10. PRIMES AND UNIQUE FACTORIZATION

**Definition 10.1.** An integer \( p \neq 0, ±1 \) is called *prime* if its only divisors are ±1, ±\( p \).

The phrase “prime number” is usually used to denote a prime positive integer. A positive integer is prime if its only positive divisors are 1 and \( p \).

The sieve of Eratosthenes: This is a method that allows one to construct rapidly a list of all primes less than a given number \( N \). We illustrate that with \( N = 50 \). One writes all the numbers from 2 to 50:

\[ \begin{align*}
2, & \ 3, \ 4, \ 5, \ 6, \ 7, \ 8, \ 9, \ 10, \ 11, \ 12, \ 13, \ 14, \ 15, \ 16, \ 17, \ 18, \ 19, \ 20, \ 21, \ 22, \ 23, \ 24, \ 25, \ 26, \ 27, \ 28, \ 29, \\
& \ 30, \ 31, \ 32, \ 33, \ 34, \ 35, \ 36, \ 37, \ 38, \ 39, \ 40, \ 41, \ 42, \ 43, \ 44, \ 45, \ 46, \ 47, \ 48, \ 49, \ 50
\end{align*}\]

The first number on the list is prime. This is 2. We write it in bold-face and cross all its multiples (we denoting crossing out by an underline):

\[ \begin{align*}
2, & \ 3, \ 4, \ 5, \ 6, \ 7, \ 8, \ 9, \ 10, \ 11, \ 12, \ 13, \ 14, \ 15, \ 16, \ 17, \ 18, \ 19, \ 20, \ 21, \ 22, \ 23, \ 24, \ 25, \ 26, \ 27, \ 28, \ 29, \\
& \ 30, \ 31, \ 32, \ 33, \ 34, \ 35, \ 36, \ 37, \ 38, \ 39, \ 40, \ 41, \ 42, \ 43, \ 44, \ 45, \ 46, \ 47, \ 48, \ 49, \ 50
\end{align*}\]

The first number on the list not in bold-face and not crossed out is prime. This is 3. We write it in bold-face and cross all its multiples (we denoting crossing out by an underline):

\[ \begin{align*}
2, & \ 3, \ 4, \ 5, \ 6, \ 7, \ 8, \ 9, \ 10, \ 11, \ 12, \ 13, \ 14, \ 15, \ 16, \ 17, \ 18, \ 19, \ 20, \ 21, \ 22, \ 23, \ 24, \ 25, \ 26, \ 27, \ 28, \ 29, \\
& \ 30, \ 31, \ 32, \ 33, \ 34, \ 35, \ 36, \ 37, \ 38, \ 39, \ 40, \ 41, \ 42, \ 43, \ 44, \ 45, \ 46, \ 47, \ 48, \ 49, \ 50
\end{align*}\]

\footnote{Eratosthenes of Cyrene, 276BC - 194BC, was a Greek mathematician who is famous for his work on prime numbers and for measuring the diameter of the earth. For more see http://www-groups.dcs.st-and.ac.uk/~history/Biographies/Eratosthenes.html}
The first number on the list not in bold-face and not crossed out is prime. This is 5. We write it in bold-face and cross all its multiples (we denoting crossing out by an underline):

\[ 2, 3, 4, 5, 6, 7, 8, 9, 10, 11, 12, 13, 14, 15, 16, 17, 18, 19, 20, 21, 22, 23, 24, 25, 26, 27, 28, 29, 30, 31, 32, 33, 34, 35, 36, 37, 38, 39, 40, 41, 42, 43, 44, 45, 46, 47, 48, 49, 50 \]

The first number on the list not in bold-face and not crossed out is prime. This is 7. We write it in bold-face and cross all its multiples (we denoting crossing out by an underline):

\[ 2, 3, 4, 5, 6, 7, 8, 9, 10, 11, 12, 13, 14, 15, 16, 17, 18, 19, 20, 21, 22, 23, 24, 25, 26, 27, 28, 29, 30, 31, 32, 33, 34, 35, 36, 37, 38, 39, 40, 41, 42, 43, 44, 45, 46, 47, 48, 49, 50 \]

The next number 11 is already greater than \(\sqrt{N} = \sqrt{50} \approx 7.071\ldots\). So we stop, because any number is a product of prime numbers (see below) and so has any number less or equal to \(N\), which is not prime, has a prime divisor smaller or equal to \(\sqrt{N}\). Thus, any number left on our list is prime.

\[ 2, 3, 4, 5, 6, 7, 8, 9, 10, 11, 12, 13, 14, 15, 16, 17, 18, 19, 20, 21, 22, 23, 24, 25, 26, 27, 28, 29, 30, 31, 32, 33, 34, 35, 36, 37, 38, 39, 40, 41, 42, 43, 44, 45, 46, 47, 48, 49, 50 \]

**Theorem 10.2.** *(The Fundamental Theorem of Arithmetic)* Every non-zero integer \(n\) is a product of primes. *(We allow the empty product, equal by definition to 1).* That is, one can write every non-zero integer \(n\) as

\[ n = \epsilon p_1 p_2 \cdots p_m, \]

where \(\epsilon = \pm 1\) and \(0 \leq p_1 \leq p_2 \leq \cdots \leq p_m\) are primes \((m \geq 0)\). Moreover, this way of writing \(n\) is unique.

**Proof.** We first show \(n\) can be written this way. We may assume \(n\) is positive (if \(n\) is negative, apply the statement for \(-n\), \(-n = p_1 p_2 \cdots p_m\) and thus \(n = -1 \cdot p_1 p_2 \cdots p_m\)).

**Lemma 10.3.** Every positive integer is a product of primes numbers. *(We allow the empty product, equal by definition to 1).*

**Proof.** Suppose not. Then the set of integers \(S\) that are not a product of prime numbers has a minimal element, say \(n_0\). \(n_0\) is not one, or a prime, because in those cases it is a product of primes. Thus, there are integers \(1 < s < n_0, 1 < t < n_0\) such that \(n_0 = st\). Note that \(s, t\) are not in \(S\) because they are smaller than \(n_0\). Thus, \(s = q_1 q_2 \cdots q_a\) is a product of primes, \(t = r_1 r_2 \cdots r_b\) is a product of primes and therefore \(n = q_1 q_2 \cdots q_a r_1 r_2 \cdots r_b\) is also a product of primes. This is a contradiction to our initial assumption that there are positive integers that are not a product of prime numbers. Thus, every positive integer is a product of prime numbers. \(\square\)
Choosing the sign $\epsilon$ appropriately and ordering the primes in increasing order we conclude that any non-zero integer $n = \epsilon p_1 p_2 \cdots p_m$, where $\epsilon = \pm 1$ and $0 \leq p_1 \leq p_2 \leq \cdots \leq p_m$ are primes. We now show uniqueness. For this we need the following important fact.

**Proposition 10.4.** Let $p$ be a positive integer. The following are equivalent: (i) $p$ is a prime number; (ii) if $p|ab$ then $p|a$ or $p|b$.

**Proof.** Suppose $p$ is prime and $p|ab$. If $p \nmid a$ then $\gcd(p, a) = 1$ and so, as we have already seen, $p|b$.

Now suppose that $p$ satisfies (ii). If $p = st$ then $p|st$ and so $p|s$, say. So $s = ps'$ and $p = ps't$. But we must have then that $s' = t = 1$, because $s, t$ are positive integers. So $p$ has no proper divisors and hence is prime. \qed

We now finish the proof of the theorem. Suppose that

$$n = \epsilon p_1 p_2 \cdots p_m,$$

and also

$$n = \mu q_1 q_2 \cdots q_t,$$

are two expressions of $n$ as in the statement of the theorem. First, $\epsilon$ is negative if and if $n$ is, and the same holds for $\mu$. So $\epsilon = \mu$. We may then assume $n$ is positive and $\epsilon = \mu = 1$ and we argue by induction on $n$. The case $n = 1$ is clear: a product of one or more primes will be greater than 1 so the only way to express $n$ is as the empty product. Assume the statement holds for $1, 2, \ldots, n - 1$ and consider two factorizations of $n$:

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

and

$$n = q_1 q_2 \cdots q_t.$$

First, note that $m \geq 1$ and $t \geq 1$ because $n > 1$. Assume that $p_1 \leq q_1$ (the argument in the other case goes the same way). We have $p_1|n$ and so $p_1|q_1 q_2 \cdots q_t$. It follows that $p_1$ divides some $q_i$ but then, $q_i$ being prime, $p_1 = q_i$. Furthermore, $p_1 \leq q_1 \leq q_i = p_1$, so $p_1 = q_1$. We then have the factorizations

$$\frac{n}{p_1} = p_2 \cdots p_m = q_2 \cdots q_t.$$

Since $n/p_1 < n$ we may apply the induction hypothesis and conclude that $m = t$ and $p_i = q_i$ for all $i$. \qed

We next derive some consequences of the fundamental theorem of arithmetic. The theorem exhibits the prime numbers as the building blocks of the integers. In itself, it doesn’t tell us if there are finitely or infinitely many such building blocks, that is, such primes.

**Theorem 10.5.** (Euclid\footnote{Euclid of Alexandria, 325BC - 265BC, was a Greek mathematician best known for his treatise on geometry: The Elements. This influenced the development of Western mathematics for more than 2000 years. For more see http://www-groups.dcs.st-and.ac.uk/%7Ehistory/Biographies/Euclid.html}) There are infinitely many prime numbers.
Proof. Let \( p_1, p_2, \ldots, p_n \) be distinct prime numbers. We show then that there is a prime not in this list. It follows that there couldn’t be only finitely many prime numbers.

Consider the integer \( n = p_1 p_2 \cdots p_n + 1 \) and its prime factorization. Let \( q \) be a prime dividing \( n \). If \( q \in \{ p_1, p_2, \ldots, p_n \} \) then \( q | p_1 p_2 \cdots p_n \) and so \( q | (n - p_1 p_2 \cdots p_n) \), that is \( q | 1 \), which is a contradiction. \( \square \)

So! We know every integer is a product of prime numbers, we know there are infinitely many prime numbers. That teaches us about the integers, and invites some more questions:

- **How frequent are the prime numbers?** The Prime Number Theorem asserts that the number of primes in the interval \([1, n]\) is roughly \( n/\log n \), in the sense that the ratio between the true number and the estimate \( n/\log n \) approaches 1 as \( n \) goes to infinity. The result was conjectured by Gauss\(^6\) at the age of 15 or 16 and proven by J. Hadamard and Ch. de la Vallée Poussin in 1896.

- **How small can the gaps between consecutive primes be?** For example, we have \((3, 5), (5, 7), (11, 13), (17, 19), \ldots\) are there infinitely many such pairs? The answer is believed to be yes but no one has proven it yet (Fall 2007). This is called the **Twin Prime Conjecture**.

- **How far does one need to go until the next prime shows up?** For example, it is known that there is always a prime between \( n \) and \( 2n \), but this is a difficult result.

- **What about adding primes?** Goldbach’s conjecture asserts that every even integer greater than 2 is the sum of two prime numbers \( 4 = 2 + 2, 6 = 3 + 3, 8 = 3 + 5, 10 = 3 + 7, 12 = 5 + 7, 14 = 3 + 11, 16 = 5 + 11, \ldots \). It has been verified (Fall 2007) up to \( n \leq 4 \times 10^{17} \), but no proof is currently known.

### 10.1. Further applications of the Fundamental Theorem of Arithmetic.

**Proposition 10.6.** Let \( a, b \) be non-zero integers. Then \( a | b \) if and only if \( a = \epsilon p_1^{a_1} \cdots p_m^{a_m} \) and \( b = \mu p_1^{a'_1} \cdots p_m^{a'_m} q_1^{b_1} \cdots q_t^{b_t} \) (products of distinct primes) with \( a'_i \geq a_i \) for all \( i = 1, \ldots, m \).

**Proof.** Clearly for such factorizations it follows that \( a | b \), in fact

\[
b/a = (\mu/\epsilon) p_1^{a'_1 - a_1} \cdots p_m^{a'_m - a_m} q_1^{b_1} \cdots q_t^{b_t}.
\]

Conversely, if \( a | b \), write \( a = \epsilon p_1^{a_1} \cdots p_m^{a_m} \) and \( b/a = \nu p_1^{a''_1} \cdots p_m^{a''_m} q_1^{b_1} \cdots q_t^{b_t} \), with \( \nu = \pm 1 \) and \( a''_i \geq 0 \). Then \( b = (\nu \epsilon) p_1^{a_1 + a''_1} \cdots p_m^{a_m + a''_m} q_1^{b_1} \cdots q_t^{b_t} \) and let \( \mu = \nu \epsilon, a'_i = a_i + a''_i \). \( \square \)

---

\(^6\)Johann Carl Friedrich Gauss, 1777 - 1855, worked in a wide variety of fields in both mathematics and physics including number theory, analysis, differential geometry, geodesy, magnetism, astronomy and optics. His work has had an immense influence in many areas. For more see [http://www-groups.dcs.st-and.ac.uk/~history/Biographies/Gauss.html](http://www-groups.dcs.st-and.ac.uk/~history/Biographies/Gauss.html)
Corollary 10.7. Let \( a = p_1^{a_1} \cdots p_m^{a_m}, b = p_1^{b_1} \cdots p_m^{b_m} \) with \( p_i \) distinct prime numbers and \( a_i, b_i \) non-negative integers. (Any two positive integers can be written this way). Then
\[
\gcd(a, b) = p_1^{\min(a_1, b_1)} \cdots p_m^{\min(a_m, b_m)}.
\]

The next proposition establishes the existence of real numbers that are not rational. It can be generalized considerably. In fact, it is known that in randomly choosing a number in the interval \([0, 1]\) the probability of picking a rational number is zero. So, though there are infinitely many rational numbers, even in the interval \([0, 1]\), they are still a rather meagre set inside the real numbers. The advantage of the following proposition is that it shows that a specific number is irrational. We will need the following result.

Proposition 10.8. Any non-zero rational number \( q \) can be written as
\[
q = \epsilon p_1^{a_1} \cdots p_m^{a_m},
\]
where \( \epsilon = \pm 1 \), the \( p_i \) are distinct prime numbers and \( a_1, \ldots, a_m \) are non-zero integers (possibly negative). Moreover, this expression is unique.

Proof. By definition, for some integers \( a, b \) we have \( q = a/b \). Let us write:
\[
a = \epsilon_a r_1^{s_1} \cdots r_n^{s_n}, \quad b = \epsilon_b r_1^{t_1} \cdots r_n^{t_n},
\]
where \( \epsilon_a, \epsilon_b \in \{\pm 1\} \), the \( r_i \) are distinct primes and \( s_i, t_i \) are non-negative integers, possibly zero. This is always possible to do because we allow zero exponents here. Then, clearly,
\[
q = (\epsilon_a/\epsilon_b)r_1^{s_1-t_1} \cdots r_n^{s_n-t_n}.
\]
Now, omitting the primes such that \( s_i-t_i = 0 \) from the list and calling the remaining primes \( p_1, \ldots, p_m \), and letting \( \epsilon = \epsilon_a/\epsilon_b \), we obtain an expression as desired.

Suppose that we have two such expressions for \( q \). Again, by allowing zero exponents, we can consider the following situation WLOG:
\[
q = \epsilon p_1^{a_1} \cdots p_m^{a_m} = \epsilon' p_1^{a'_1} \cdots p_m^{a'_m}.
\]
Since \( \epsilon, \epsilon' \) determine the sign of \( q \), they must be equal and we need to show that \( a_i = a'_i \) for all \( i \). Dividing through, we get an expression of the form,
\[
1 = p_1^{c_1} \cdots p_m^{c_m},
\]
where \( c_i = a_i - a'_i \) and we need to show all the \( c_i \) are zero. WLOG, by rearranging the primes, we may assume that \( c_1, \ldots, c_t \) are negative and \( c_{t+1}, \ldots, c_m \) are non-negative. We conclude that
\[
p_1^{-c_1} \cdots p_t^{-c_t} = p_1^{c_{t+1}} \cdots p_m^{c_m}.
\]
In this expression there are no negative exponents. Thus, from unique factorization for integers, since all the primes are distinct all the powers must be zero. \( \square \)

Proposition 10.9. \( \sqrt{2} \) is not a rational number.
\textit{Proof.} Suppose it is, and write $\sqrt{2} = p_1^{a_1} \cdots p_m^{a_m}$, distinct primes with non-zero exponents (possibly negative). Then

$$2 = p_1^{2a_1} \cdots p_m^{2a_m}$$

must be the unique factorization of 2. However, 2 is prime. So there must be only one prime on the right hand side, i.e. $m = 1$. Then $2 = p_1^{2a_1}$ and we must have $p_1 = 2$ and $2a_1 = 1$. But this contradicts the fact that $a_1$ is an integer.

Here is another proof:

Suppose that $\sqrt{2}$ is rational and write $\sqrt{2} = m/n$, where $(m, n) = 1$. Then

$$2n^2 = m^2.$$ 

This implies that $2|m^2$. If $m = 2k + 1$ then $m^2 = 4k^2 + 4k + 1$. Since $2|(4k^2 + 4k)$ it would follow that $2|1$ which is a contradiction. (In simple language, what we are saying is that $m$ is odd implies that $m^2$ is odd.) Thus, $m$ is even. Then $2|m$. Say $m = 2k$. It follows that $2n^2 = 4k^2$ and so $n^2 = 2k^2$. Therefore, $2|n^2$ and so $2|n$. This means that 2 divides both $n$ and $m$, contrary to our assumption. Thus, assuming $\sqrt{2}$ is rational leads to a contradiction and so $\sqrt{2}$ is not a rational number. □

Though the second proof is more elementary, the first proof is better. It is much easier to generalize and indeed in a similar way, one can show $\sqrt{3}$, $\sqrt{17}$ etc. are irrational.

It is also known that $e$ is irrational (not too hard) and $\pi$ is irrational (hard). But it is still an open question (Fall 2007) if Euler’s constant

$$\gamma = \lim_{n \to \infty} \left(1 + \frac{1}{2} + \frac{1}{3} + \cdots + \frac{1}{n} - \log(n)\right) \approx 0.57721$$

is rational or not (it is believed to be irrational; if $\gamma$ is rational, it was proved that its denominator has to have more than $10^{242080}$ digits!).
11. Relations

A relation on a set $S$ is best described as a subset $\Gamma \subset S \times S$. For each $s \in S$, $s$ is related to $t$ if $(s, t) \in \Gamma$. Though the format reminds one of functions, the actual relevance of the notion of functions is minimal. For example, usually for a given $s$ there will be many elements $t$ such that $(s, t) \in \Gamma$, which is the opposite of what we have for functions, where there is only one $t$. We shall usually denote that $x$ is related to $y$, namely that $(x, y) \in \Gamma$, by $x \sim y$.

Note that so far the definition is wide enough to allow any $\Gamma$. A relation is called transitive if $x \sim y$ and $y \sim z$ implies $x \sim z$. A relation is called a partial order if it is transitive and in addition we don’t have both $x \sim y$ and $y \sim x$. We then use the notation $x < y$ for $x \sim y$. There may very well be $x, y$ for which neither $x < y$ nor $y < x$ holds. A linear order (or a simple order) is a partial order such that for every $x, y$ we have either $x < y$ or $y < x$.

Another important class of relations, even more important for this course, are the equivalence relations. They are very far from order relations. A relation is called an equivalence relation if it satisfies the following properties:

1. (Reflexive) For every $x$ we have $x \sim x$.
2. (Symmetric) If $x \sim y$ then $y \sim x$.
3. (Transitive) If $x \sim y$ and $y \sim z$ then $x \sim z$.

Equivalence relations arise when one wishes to identify elements in a given set $S$ according to some principle. Implicit in the word “identify” is that $x$ is identified with $x$, if $x$ is identified with $y$ then $y$ is identified with $x$, and that if $x$ is identified with $y$ and $y$ is identified with $z$ then, by all accounts, $x$ should be identified with $z$ too. That is, we have an equivalence relation.

**Lemma 11.1.** Let $\sim$ be an equivalence relation on a set $S$. Define the equivalence class $[x]$ of an element $x \in S$ as follows:

$$[x] = \{ y : y \in S, x \sim y \}.$$  

This is a subset of $S$. The following holds:

1. Two equivalence classes are either disjoint or equal.
2. $S$ is a disjoint union of equivalent classes.

Conversely, if $S$ is a disjoint union $S = \bigcup_{i \in I} U_i$ of non-empty sets (this is called a partition of $S$) then there is a unique equivalence relation on $S$ for which the $U_i$ are the equivalence classes.

**Proof.** Let $x, y$ be elements of $S$ and suppose that $[x] \cap [y] \neq \emptyset$. Then, there is an element $z$ such that $x \sim z, y \sim z$. Since $\sim$ is symmetric also $z \sim y$ and using transitivity $x \sim y$. Now, if $s \in [y]$ then $y \sim s$ and by transitivity $x \sim s$ and so $s \in [x]$ and we showed $[y] \subset [x]$. Since $x \sim y$ also $y \sim x$ and the same argument gives $[x] \subset [y]$. We conclude that $[x] = [y]$.
Every element of \( S \) lies in the equivalence class of itself. It follows that \( S \) is a disjoint union of equivalence classes.

To prove the second part of the lemma, we define that \( x \sim y \) if both \( x \) and \( y \) lie in the same set \( U_i \). It is clearly reflexive and symmetric. It is also transitive: \( x \sim y \) means \( x, y \in U_i \) for some \( i \), \( y \sim z \) means \( y, z \in U_j \) for some \( j \). But there is a unique \( U_i \) containing \( y \) because the union is a disjoint union. That is \( U_i = U_j \) and so \( x, z \in U_j \), meaning \( x \sim z \). The equivalence classes are clearly the \( U_i \). \( \square \)

We introduce the following terminology: we say that a set \( \{ x_i : i \in I \} \), \( I \) some index set, is a complete set of representatives if the equivalence classes \( [x_i] \) are disjoint and \( S = \bigcup_{i \in I} [x_i] \). This means that every equivalence class if of the form \( [x_i] \) for some \( i \in I \) and no two equivalence classes \( [x_i], [x_j] \) for \( i, j \in I \) are the same, unless \( x_i = x_j \).

12. Congruence relations

Let \( n \) be a positive integer. Define a relation \( x \sim y \) if \( n | (x - y) \) (we shall also write that as \( x \equiv y \pmod{n} \), or simply \( x \equiv y \) if \( n \) is clear from the context). We say that \( x \) is congruent to \( y \) modulo \( n \).

**Lemma 12.1.** Congruence modulo \( n \) is an equivalence relation on \( \mathbb{Z} \). The set \( \{0, 1, \ldots, n-1\} \) is a complete set of representatives.

**Proof.** First \( n | (x - x) \) so \( x \equiv x \) and the relation is reflexive. If \( n | (x - y) \) then \( n | (x - y) = y - x \), so the relation is symmetric. Suppose \( n | (x - y), n | (y - z) \) then \( n | (x - y) + (y - z) = x - z \) and so the relation is transitive too.

Let \( x \) be any integer and write \( x = qn + r \) with \( 0 \leq r < n \). Then \( x - r = qn \) and so \( x \equiv r \). It follows that every equivalence class is represented by some \( r \in \{0, 1, \ldots, n-1\} \). The equivalence classes defined by elements of \( \{0, 1, \ldots, n-1\} \) are disjoint. If not, then for some \( 0 \leq i < j < n \) we have \( i \equiv j \), that is, \( n | (j - i) \). But \( 0 < j - i < n \) and we get a contradiction. \( \square \)

**Theorem 12.2.** Denote the equivalence classes of congruence modulo \( n \) by \( \overline{0}, \overline{1}, \ldots, \overline{n-1} \). Denote this set by \( \mathbb{Z}/n\mathbb{Z} \). The set \( \mathbb{Z}/n\mathbb{Z} \) is a commutative ring under the following operations:

\[
\overline{i} + \overline{j} = \overline{i + j}, \quad \overline{i} \cdot \overline{j} = \overline{ij}.
\]

The neutral element for addition is \( \overline{0} \), for multiplication \( \overline{1} \) and the inverse of \( \overline{i} \) is \( \overline{-i} = \overline{n-i} \).

Before proving the theorem we illustrate the definitions in a numerical example:

**Example 12.3.** We take \( n = 13 \) and calculate \( \overline{5} \cdot \overline{6} = \overline{5} \). First, \( \overline{5} \cdot \overline{6} = \overline{30} = \overline{4} \). Then \( \overline{4} - \overline{5} = \overline{4} + \overline{-5} = \overline{4-5} = \overline{-1} = \overline{12} \). Note that we could have also calculated \( \overline{4} - \overline{5} = \overline{4} + \overline{8} = \overline{12} \), or \( \overline{5} \cdot \overline{-5} = \overline{5(6-1)} = \overline{5} \cdot \overline{5} = \overline{25} = \overline{12} \).
Modular arithmetic, that is calculating in the ring $\mathbb{Z}/n\mathbb{Z}$, is sometimes called “clock arithmetic”. The reason is the following. The usual clock is really displaying hours modulo 12. When 5 hours pass from the time 10 o’clock the clock shows 3. Note that $3 \equiv 15 \pmod{12}$. We are used to adding hours modulo 12 (or modulo 24, for that matter), but we are not used to multiplying hours, that doesn’t make sense. However, if you’d like you can think about multiplication as repeated addition $5 \cdot 3 = 5 + 5 + 5$. So, in that sense, we are already familiar with the operations modulo 12 and the definitions above are a generalization.

Continuing with our numerical example, let us try and solve the equation $4x + 2 = 7$ in $\mathbb{Z}/13\mathbb{Z}$. Now, and from now on, we are just writing $4, 2, 7$ etc. for $\bar{4}, \bar{2}, \bar{7}$. So we need to solve $4x = 5$. We are now looking for a residue class $r$ modulo 13 so that $4r \equiv 1 \pmod{13}$. We guess that $r = 10$ and check: $4 \cdot 10 = 40 \equiv 1 \pmod{13}$. Then, $4x = 5$ implies $10 \cdot 4x \equiv x \equiv 50 \equiv 11 \pmod{13}$. Thus, the only possibility is $x = 11$. We go back to the original equation $4x = 5$ and verify that $4 \cdot 11 \equiv 5 \pmod{13}$. We found the solution $x = 11$. We remark that in general such an $r$ need not exist if the modulos $n$ is not a prime, and that in that case one may need to go back to the original equation and verify that indeed the solution to the reduced equation solves the original equation. These issues will be discussed later.

**Example 12.4.** As another example, we give that addition and multiplication table of $\mathbb{Z}/5\mathbb{Z}$.

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**Proof.** (Of Theorem) We first prove that the operations do not depend on the representatives for the equivalence classes that we have chosen. Suppose $\bar{i} = \bar{i}', \bar{j} = \bar{j}'$, where $i, i', j, j'$ need not be in the set $\{0, 1, 2, \ldots, n - 1\}$. We defined $\bar{i} + \bar{j} = \bar{i} + \bar{j}$. We need to check that this is the same as $\bar{i}' + \bar{j}'$. Since $\bar{i} = \bar{i}', n|(i - i')$ and similarly $n|(j - j')$. Therefore, $n|(i + j) - (i' + j')$, that is, $\bar{i} + \bar{j} = \bar{i}' + \bar{j}'$.

We also need to show that $\bar{ij} = \bar{i'j'}$. But, $ij - i'j' = ij - ij' + ij' - i'j' = i(j - j') + j'(i - i')$ and so $n|(ij - i'j')$.

The verification of the axioms is now easy if we make use of the fact that $\mathbb{Z}$ is a commutative ring:

1. $\bar{i} + \bar{j} = \bar{i} + \bar{j} = \bar{j} + \bar{i}$.
2. $(\bar{i} + \bar{j}) + \bar{k} = \bar{i} + \bar{j} + \bar{k} = \bar{i} + (\bar{j} + \bar{k})$. Note that at this point we used the simplification that we can use any representative of the equivalence class to carry out the operations. Had we insisted on always using the representative in the set $\{0, 1, 2, \ldots, n - 1\}$ we would usually need to replace $i + j$ by its representative in that set and things would be turning messy. Now, $\bar{i} + \bar{j} + \bar{k} = \bar{i} + (\bar{j} + \bar{k}) = \bar{i} + \bar{j} + \bar{k}$. 

(3) \( \overline{0} + \overline{i} = \overline{0+i} = \overline{i} \).

(4) \( \overline{i} + \overline{-i} = \overline{i + (-i)} = \overline{0} \). Note that \( \overline{-i} = \overline{n-i} \).

(5) \( (\overline{i} \cdot \overline{j})\overline{k} = \overline{ij} \cdot \overline{k} = (\overline{ij})\overline{k} = \overline{i(jk)} = \overline{i} \cdot \overline{jk} = \overline{i} (\overline{j} \cdot \overline{k}) \).

(6) \( 1 \cdot \overline{i} = \overline{1} \cdot \overline{i} = \overline{i} \cdot \overline{1} = \overline{i} \cdot 1 = \overline{i}. \)

(7) \( \overline{i}(\overline{j} + \overline{k}) = \overline{i} \cdot \overline{j} + \overline{i} \cdot \overline{k} = \overline{ij} + \overline{ik} = \overline{ij} + \overline{ik} = \overline{i} \cdot \overline{j} + \overline{i} \cdot \overline{k}. \) Similarly, \( \overline{(j + k)i} = \overline{j + k} \cdot \overline{i} = (\overline{j} + \overline{k})\overline{i} = \overline{ji} + \overline{ki} = \overline{ji} + \overline{ki} = \overline{j} \cdot \overline{i} + \overline{k} \cdot \overline{i}. \)

Furthermore, this is a commutative ring: \( \overline{i} \cdot \overline{k} = \overline{ik} = \overline{ki} = \overline{k} \cdot \overline{i} \). \( \square \)

In the proof we saw that the ring properties of \( \mathbb{Z}/n\mathbb{Z} \), the set of equivalence classes modulo \( n \), all follow from the ring properties of \( \mathbb{Z} \). We shall later see that this can be generalized to any ring \( R \): if we impose a correct notion of an equivalence relation, the equivalence classes themselves will form a ring and the fact that the axioms hold follows from the fact they hold for \( R \).

**Theorem 12.5.** \( \mathbb{Z}/n\mathbb{Z} \) is a field if and only if \( n \) is prime.

Before providing the proof we introduce some terminology. Let \( R \) be a ring, \( x \in R \) a non-zero element. \( x \) is called a zero divisor if there is an element \( y \neq 0 \) such that either \( xy = 0 \) or \( yx = 0 \) (or both).

**Lemma 12.6.** Let \( R \) be a commutative ring. If \( R \) has zero divisors then \( R \) is not a field.

**Proof.** Let \( x \neq 0 \) be a zero divisor and let \( y \neq 0 \) be an element such that \( xy = 0 \). If \( R \) is a field then there is an element \( z \in R \) such that \( zx = 1 \). But then \( z(xy) = z \cdot 0 = 0 \) and also \( z(xy) = (zx)y = 1 \cdot y = y \) so \( y = 0 \) and that is a contradiction. \( \square \)

**Proof.** (Of Theorem) If \( n = 1 \) then \( \mathbb{Z}/n\mathbb{Z} \) has a single element and so \( 0 = 1 \) in that ring. Therefore, it is not a field. Suppose that \( n > 1 \) and \( n \) is not prime, \( n = ab \) where \( 1 < a < n, 1 < b < n \). Then \( \overline{a} \neq 0, \overline{b} \neq 0 \) but \( \overline{a} \cdot \overline{b} = \overline{ab} = \overline{n} = \overline{0} \). So \( \mathbb{Z}/n\mathbb{Z} \) has zero divisors and thus is not a field.

Suppose now that \( n \) is prime and let \( \overline{a} \neq 0 \). That is, \( n \nmid a \), which, since \( n \) is prime, means that \( (n,a) = 1 \). Consider the list of elements

\[ 0 \cdot \overline{a}, 1 \cdot \overline{a}, \ldots, n-1 \cdot \overline{a}. \]

We claim that they are distinct elements of \( \mathbb{Z}/n\mathbb{Z} \). Suppose that \( \overline{i} \cdot \overline{a} = \overline{j} \cdot \overline{a} \), for some \( 0 \leq i \leq j \leq n-1 \) then \( \overline{ia} = \overline{ja} \), which means that \( n|(ia-ja) = (i-j)a \). Since \( (n,a) = 1 \), it follows that \( n|(i-j) \) but that means \( i = j \). Thus, the list \( 0 \cdot \overline{a}, 1 \cdot \overline{a}, \ldots, n-1 \cdot \overline{a} \) contains \( n \) distinct elements of \( \mathbb{Z}/n\mathbb{Z} \) and so it must contain \( 1 \). That is, there’s an \( i \) such that \( \overline{i} \cdot \overline{a} = \overline{1} \) and therefore \( \overline{a} \) is invertible. \( \square \)

Let \( p \) be a prime number. We denote \( \mathbb{Z}/p\mathbb{Z} \) also by \( \mathbb{F}_p \). It is a field with \( p \) elements. In fact, any finite field, that is any field with finitely many elements, has cardinality a power of a prime and for any prime power there is a field with that cardinality. Finite fields, such as \( \mathbb{F}_p \), play an important role in coding and cryptography as well as in pure mathematics.
12.1. Fermat’s little theorem.

**Theorem 12.7.** (Fermat’s little theorem) Let $p$ be a prime number. Let $a \not\equiv 0 \pmod{p}$ then

\[ a^{p-1} \equiv 1 \pmod{p}. \]

Before proving the theorem we prove state two auxiliary statements whose proofs are delegated to the assignments.

**Lemma 12.8.** We have $p \mid \binom{p}{i}$ for every $1 \leq i \leq p - 1$.

**Lemma 12.9.** Let $R$ be a commutative ring and $x, y \in R$. Interpret $\binom{n}{i}$ as adding $\binom{n}{i}$ times the element 1. Then the binomial formula holds in $R$:

\[ (x + y)^n = \sum_{i=0}^{n} \binom{n}{i} x^i y^{n-i}. \]

**Proof.** (Of Fermat’s little theorem) We prove that by induction on $1 \leq a \leq p - 1$. For $a = 1$ the result is clear. Suppose the result for $a$ and consider $a + 1$, provided $a + 1 < p$. We have, by the binomial formula,

\[ (a + 1)^p = \sum_{i=0}^{p} \binom{p}{i} a^i \]
\[ = 1 + \binom{p}{1} a + \binom{p}{2} a^2 + \cdots + \binom{p}{p-1} a^{p-1} + a^p \]
\[ = 1 + a^p \quad \text{(using the lemma)} \]
\[ = 1 + a \quad \text{(using the induction hypothesis)} \]

Since $1 + a \not\equiv 0 \pmod{p}$ it has an inverse $y$ in $\mathbb{F}_p$, $y(1 + a) \equiv 1$ and we get $y(1 + a)^p \equiv y(1 + a)$, that is $(1 + a)^{p-1} \equiv 1$. □

**Example 12.10.** We calculate $2^{100}$ modulo 13. We have $2^{100} = 2^{96}2^4 = (2^{12})^82^4 \equiv 2^4 \equiv 3 \pmod{13}$.

Fermat’s little theorem gives a criterion for numbers to be composite. Let $n$ be a positive integer. If there is $1 \leq a \leq n - 1$ such that $a^{n-1} \not\equiv 1 \pmod{n}$ then $n$ is not prime. Unfortunately, it is possible that for every $1 \leq a \leq n - 1$ such that $(a,n) = 1$, one has $a^{n-1} \equiv 1 \pmod{n}$ and yet $n$ is not prime. Thus, this test fails to recognize such $n$ as composite numbers. Such numbers are called Carmichael numbers. There are infinitely many such numbers. The first being

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7Pierre de Fermat, 1601 - 1665, was a French lawyer and government official most remembered for his work in number theory; in particular for Fermat’s Last Theorem. He is also important in the foundations of the calculus. For more, see http://www-groups.dcs.st-and.ac.uk/~history/Biographies/Fermat.html
Primality testing programs first test divisibility by small primes available to the program as pre-computed data and then choose randomly some \( 1 \leq a < n \): if \((a, n) \neq 1\) then \(n\) is not prime. If \((a, n) = 1\) the program calculated \(a^{n-1} \pmod{n}\). If the result is not \(1 \pmod{n}\) then \(n\) is not prime. If the result is \(1\), the program chooses another \(a\). After a certain number of tests, say 10, if \(n\) passed all the tests it is declared as “prime”, though there is no absolute reassurance it is indeed a prime. We remark that calculating \(a^{n-1} \pmod{n}\) can be done quickly. One calculates \(a, a^2, a^4, a^8, a^{16}, \ldots \) modulo \(n\), as long as the power is less than \(n\). This can be done rapidly. One then expresses \(n\) in base 2 to find the result. Here is an example: Let us calculate \(3^{54} \pmod{55}\) (random choice of numbers). We have \(3, 3^2 = 9, 3^4 = 81 = 26, 3^8 = 676 = 16, 3^{16} = 16^2 = 256 = 36, 3^{32} = 36^2 = 1296 = 31.\) Now, \(54 = 2 + 4 + 16 + 32\) and so \(3^{54} = 9 \cdot 26 \cdot 36 \cdot 31 = 4.\) In particular, 55 is not a prime – not that I’m particularly proud in drawing this conclusion...

It is important to note that there is a polynomial time algorithm to decide, without any doubt, if an integer is prime. Such an algorithm was discovered by Agrawal, Kayal and Saxena in 2002. It is important to note that the algorithm does not produce a decomposition of \(n\) in case \(n\) is composite. Such an algorithm will compromise the very backbone of e-commerce and military security. For more, see [http://www.ams.org/notices/200305/fea-bornemann.pdf](http://www.ams.org/notices/200305/fea-bornemann.pdf)

### 12.2. Solving equations in \(\mathbb{Z}/n\mathbb{Z}\)

There is no general method to solving polynomials equations in \(\mathbb{Z}/n\mathbb{Z}\). We just present some selected topics.

#### 12.2.1. Linear equations

We want to consider the equation \(ax + b \in \mathbb{Z}/n\mathbb{Z}\). Let us assume that \((a, n) = 1\). Then, there are integers \(u, v\) such that \(1 = ua + vn\). We remark that \(u, v\) are found by the Euclidean algorithm. Note that this implies that \(ua \equiv 1 \pmod{n}\). Thus, if \(x\) solves \(ax + b = 0\) in \(\mathbb{Z}/n\mathbb{Z}\) then \(x\) solves the equation \(uax + ub = 0 \pmod{n}\), that is \(x + ub = 0\) and so \(x = -ub\) in \(\mathbb{Z}/n\mathbb{Z}\). Conversely, if \(x = -ub\) in \(\mathbb{Z}/n\mathbb{Z}\) where \(ua = 1\) in \(\mathbb{Z}/n\mathbb{Z}\) then \(ax = a(-ub) = -aub = -b\) in \(\mathbb{Z}/n\mathbb{Z}\).

We summarize: if \((a, n) = 1\) then the equation

\[
ax + b = 0 \pmod{n},
\]

has a unique solution \(x = -ub\), where \(u\) is such that \(ua = 1 \pmod{n}\).

Here is a numerical example: Let us solve \(12x + 3 = 0 \pmod{17}\). First \(17 = 12 + 5, 12 = 2 \cdot 5 + 2, 5 = 2 \cdot 2 + 1\), so \((12, 17) = 1\). Moreover, \(1 = 5 - 2 \cdot 2 = 5 - 2 \cdot (12 - 2 \cdot 5) = 5 \cdot 5 - 2 \cdot 12 = 5 \cdot (17 - 12) - 2 \cdot 12 = 5 \cdot 17 - 7 \cdot 12\). We see that \(-7 \cdot 12 = 1 \pmod{17}\). Thus, the solution is \(x = 7 \cdot 3 = 21 = 4 \pmod{17}\).

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8For more see [http://mathworld.wolfram.com/CarmichaelNumber.html](http://mathworld.wolfram.com/CarmichaelNumber.html)
12.2.2. **Quadratic equations.** Consider the equation \( ax^2 + bx + c = 0 \) in \( \mathbb{Z}/n\mathbb{Z} \) and assume \( n \) is a prime greater than 2. In that case, assuming that \( a \neq 0 \) modulo \( n \), there is an element \((2a)^{-1}\). One can prove that the equation has a solution if and only if \( b^2 - 4ac \) is a square in \( \mathbb{Z}/n\mathbb{Z} \) (which may or may not be the case). In case it is a square, the solutions of this equation are given by the usual formula:

\[
(2a)^{-1}(-b \pm \sqrt{b^2 - 4ac}).
\]

For example, the equation \( x^2 + x + 1 \) has no solution in \( \mathbb{Z}/5\mathbb{Z} \) because the discriminant \( b^2 - 4ac \) is in this case \( 1^2 - 4 = -3 \) in \( \mathbb{Z}/5\mathbb{Z} \) and 2 can be checked not to be a square in \( \mathbb{Z}/5\mathbb{Z} \) (one just tries: \( 0^2 = 0, 1^2 = 1, 2^2 = 4, 3^2 = 4, 4^2 = 1 \) in \( \mathbb{Z}/5\mathbb{Z} \)). One the other hand, \( x^2 + x + 1 \) can be solved in \( \mathbb{Z}/7\mathbb{Z} \). The solutions are \( 4(-1 \pm \sqrt{-3}) = -4 \pm 4\sqrt{3} = -4 \pm 8 = -4 \pm 1 = \{2, 4\} \).

When \( n \) is not prime, we shall not study the problem in this course, beyond remarking that one can proceed by trying all possibilities if \( n \) is small and that the number of solutions can be very large. For example: consider the equation \( x^3 - x \) in \( \mathbb{Z}/8\mathbb{Z} \). We can verify that its solutions are \( 0, 1, 3, 5, 7 \). There are 5 solutions but the equation has degree three. We shall later see that in any field a polynomial equation of degree \( n \) has at most \( n \) roots.

12.3. **Public key cryptography; RSA method.**

We cannot go here too much into the cryptographical practical aspects. Suffices to say that in many cryptographical applications two parties \( X \) and \( Y \) wish to exchange a secret. Given any large integer \( n \) that secret can be represented as a number modulo \( n \), and we leave it to the reader’s imagination to devise methods for that. The method proceeds as follows:

- \( X \) chooses two large primes \( p < q \).
- \( X \) calculates \( n = pq \).
- \( X \) calculates \( k = (p - 1)(q - 1) \).
- \( X \) chooses an integer \( d \) such that \( (d, k) = 1 \).
- \( X \) finds \( e \) such that \( ed \equiv 1 \pmod{k} \).
- \( X \) publishes for anyone to see the data \( e, n \). This is called the **public key**.

The rest of the data \( p, q, k, d \) is kept secret. In fact, \( p, q, k \) can be destroyed altogether and only \( d \) be kept, and kept secret. This is called the **private key**.

\( Y \), wishing to send a secret, writes it as a number \( b \) modulo \( n \), which is also relatively prime to \( n \), and sends \( b^e \pmod{n} \) to \( X \), allowing anyone interested to see that message. The point is, and this is called the **discrete log problem**, that it is very difficult to find what \( b \) is, even when one knows \( b^e \) and \( n \). Thus, someone seeing \( Y \)'s message cannot find the secret \( b \) from it.

\( X \), upon receiving \( Y \)'s message \( b^e \), calculates \((b^e)^d \).

\(^9\)The RSA method described above is named after Ron Rivest, Adi Shamir and Len Adleman, who discovered it in 1977.
Lemma 12.11. We have $b^{ed} \equiv b \pmod{n}$.

Proof. We need to show that $b^{ed} \equiv b \pmod{p}$ and $b^{ed} \equiv b \pmod{q}$. Then $p|(b^{ed} - b)$ and $q|(b^{ed} - b)$ and so (using the $p, q$ are primes and distinct), $n = pq|(b^{ed} - b)$.

The argument being symmetric, we just show $b^{ed} \equiv b \pmod{p}$. We have modulo $p$

$$b^{ed} = b^{1+vk}$$

$$= b \cdot ((b^{p-1})^{q-1})^v$$

$$= b \cdot (1^{q-1})^v \quad \text{Fermat’s little theorem}$$

$$= b.$$  

\[ \square \]

We have shown that X can retrieve Y’s secret.

Here is a numerical example:

\[ p = 10007, \quad q = 10009; \]
\[ n = p*q = 100160063 \]
\[ k = (p-1)*(q-1) = 100140048 \]
\[ d = 10001 \]
\[ e = 88695185 \]
\[ b = 3 \]
\[ b^e = 33265563 \]
\[ 33265563^{10001} = 3 \pmod{n}. \]
Part 4. Polynomials and their arithmetic

13. The ring of polynomials

Let \( R \) be a commutative ring. A good example to keep in mind is \( R = \mathbb{Z} \) or \( R = \mathbb{C} \), but our discussion allows any ring. We define the ring of polynomials over \( R \) as

\[
R[x] = \{a_n x^n + \cdots + a_1 x + a_0 : a_i \in R\}.
\]

In the definition \( n \) is any non-negative integer. Note that we allow some, or even all, coefficients to be zero. The zero polynomial 0 is the choice \( n = 0 \) and \( a_0 = 0 \). We define addition as (assume \( n \geq m \))

\[
(a_n x^n + \cdots + a_1 x + a_0) + (b_m x^m + \cdots + b_1 x + b_0) = a_n x^n + \cdots + (a_m + b_m) x^m + \cdots (a_1 + b_1) x + (a_0 + b_0).
\]

We also define multiplication by

\[
(a_n x^n + \cdots + a_1 x + a_0)(b_m x^m + \cdots + b_1 x + b_0) = c_{n+m} x^{n+m} + \cdots + c_1 x + c_0,
\]

where

\[
c_i = a_0 b_i + a_1 b_{i-1} + \cdots + a_i b_1 + a_i b_0.
\]

Note that in the formula for \( c_i \) it is entirely possible that some \( a_j \) or \( b_j \) are not defined; this happens if \( j > n \) or \( j > m \), respectively. In this case we understand \( a_j \), or \( b_j \), as zero.

Example 13.1. Take \( R = \mathbb{Z} \) then

\[
(2x^2 - x - 2) + (x^3 + x - 1) = x^3 + 2x^2 + 2x - 3,
\]

\[
(2x^2 - x - 2)(x^3 + x - 1) = 2x^5 + x^4 - x^2 - 3x + 2.
\]

A polynomial \( f(x) = a_n x^n + \cdots + a_1 x + a_0 \) is called monic if \( a_n = 1 \). It is called of degree \( n \) if \( a_n \neq 0 \). If \( f \) has degree 0, that is \( f(x) = a, a \in R, a \neq 0 \), then \( f \) is called a constant polynomial. The degree of the zero polynomial is not defined.

Proposition 13.2. With the operations defined above \( R[x] \) is a commutative ring, with zero being the zero polynomial and 1 being the constant polynomial 1. The additive inverse of \( a_n x^n + \cdots + a_1 x + a_0 \) is \(-a_n x^n - \cdots - a_1 x - a_0 \).

Since the proof is straightforward we leave it as an exercise.

Proposition 13.3. If \( R \) is an integral domain then \( R[x] \) is an integral domain. If \( f(x), g(x) \in R[x] \) are non-zero polynomials,

\[
\deg(f(x)g(x)) = \deg(f(x)) + \deg(g(x)).
\]

Proof. Say \( \deg(f(x)) = n \), \( \deg(g(x)) = m \), then by definition \( f(x) = a_n x^n + \cdots + a_1 x + a_0 \) with \( a_n \neq 0 \) and \( g(x) = b_m x^m + \cdots + b_1 x + b_0 \) with \( b_m \neq 0 \). Then \( f(x)g(x) = a_n b_m x^{n+m} + (a_n b_{m-1} + a_{n-1} b_m)x^{n+m-1} + \cdots \). Since \( R \) is an integral domain \( a_n b_m \neq 0 \) and so \( f(x)g(x) \neq 0 \) and \( \deg(f(x)g(x)) = n + m \). \( \square \)
14. Division with residue

Let \( \mathbb{F} \) be a field. We have defined the ring of polynomials \( \mathbb{F}[x] \); it is an integral domain (but is never a field; for example \( x \) does not have an inverse with respect to multiplication).

**Theorem 14.1.** Let \( f(x), g(x) \) be two polynomials in \( \mathbb{F}[x], g(x) \neq 0 \). Then, there exist unique polynomials \( q(x), r(x) \) in \( \mathbb{F}[x] \) such that

\[
  f(x) = q(x)g(x) + r(x), \quad r(x) = 0 \text{ or } \deg(r(x)) < \deg g(x).
\]

**Proof.** We first show the existence and later the uniqueness. Consider the set

\[
  S = \{ f(x) - q(x)g(x) : q(x) \in \mathbb{F}[x] \}.
\]

If \( 0 \in S \) then there is a \( q(x) \) such that \( f(x) = q(x)g(x) \) and we take \( r(x) = 0 \). Else, choose an element \( r(x) \) in \( S \) of minimal degree. Since \( r(x) \) is in \( S \) we can write \( r(x) = f(x) - q(x)g(x) \) for some \( q(x) \).

**Claim.** \( \deg(r(x)) < \deg(g(x)). \)

Let us write \( r(x) = r_n x^n + \cdots + r_1 x + r_0 \) and \( g(x) = g_m x^m + \cdots + g_1 x + g_0 \), with \( r_n \neq 0, g_m \neq 0 \).

Assume, by contradiction, that \( n \geq m \). Then \( r_1(x) = r(x) - r_n g_m^{-1} x^{n-m} g(x) = (r_{n-1} - r_n g_m^{-1} g_{m-1}) x^{n-1} + \cdots \) has degree smaller then \( r(x) \). On the other hand, \( r_1(x) = r(x) - r_n g_m^{-1} x^{n-m} g(x) = f(x) - q(x)g(x) - r_n g_m^{-1} x^{n-m} g(x) = f(x) - (q(x) + r_n g_m^{-1} x^{n-m})g(x) \) shows that \( r_1(x) \in S \). Contradiction. We have therefore established the existence of \( q(x), r(x) \) such that

\[
  f(x) = q(x)g(x) + r(x), \quad r(x) = 0 \text{ or } \deg(r(x)) < \deg g(x).
\]

We now prove uniqueness. Suppose that also

\[
  f(x) = q_1(x)g(x) + r_1(x), \quad r_1(x) = 0 \text{ or } \deg(r_1(x)) < \deg g(x).
\]

We need to show that \( q(x) = q_1(x), r(x) = r_1(x) \). We have,

\[
  (q(x) - q_1(x))g(x) = r_1(x) - r(x).
\]

The right hand side is either zero or has degree less that \( g(x) \). If it’s zero then, since \( \mathbb{F}[x] \) is an integral domain, we also have \( q(x) = q_1(x) \). If \( r(x) \neq r_1(x) \) then also \( q(x) \neq q_1(x) \) but then the degree of the left hand side is \( \deg(q(x) - q_1(x)) + \deg(g(x)) \geq \deg(g(x)) \) and we get a contradiction. \( \square \)

15. Arithmetic in \( \mathbb{F}[x] \)

In this section \( \mathbb{F} \) is a field. We denote by \( \mathbb{F}^x \) the set of non-zero elements of \( \mathbb{F} \).
15.1. Some remarks about divisibility in a commutative ring $T$. The definitions we made in §9 can be made in general and the same basic properties hold. Let $T$ be a commutative ring and $a, b \in T$. We say that $a$ divides $b$ if $b = ac$ for some $c \in T$. We have the following properties:

1. $a | b \Rightarrow a | b - a$.
2. $a | b \Rightarrow a | bd$ for any $d \in T$.
3. $a | b, a | d \Rightarrow a | (b \pm d)$.

In particular, these definition and properties hold for the ring of polynomials $R[x]$, where $R$ is a commutative ring.

15.2. GCD of polynomials.

**Definition 15.1.** Let $f(x), g(x) \in F[x]$, not both zero. The greatest common divisor of $f(x)$ and $g(x)$, denoted $\gcd(f(x), g(x))$ of just $(f(x), g(x))$, is the monic polynomial of largest degree dividing both $f(x)$ and $g(x)$. (We shall see below that there is a unique such polynomial.)

**Theorem 15.2.** Let $f(x), g(x)$ be polynomials, not both zero. The gcd of $f(x)$ and $g(x)$, $h(x) = (f(x), g(x))$, is unique and can be expressed as

$$h(x) = u(x)f(x) + v(x)g(x), \quad u(x), v(x) \in F[x].$$

It is the monic polynomial of minimal degree having such an expression. If $t(x)$ divides both $g(x)$ and $f(x)$ then $t(x)|h(x)$.

**Proof.** Consider the following set of monic polynomials

$$S = \{a(x) : a(x) = u(x)f(x) + v(x)g(x) \text{ for some } u(x), v(x) \in F[x], a(x) \text{ monic}\}.$$ 

It is a non-empty set because if $f(x) = bx^n + \text{l.o.t.}$ then $b^{-1}f(x) \in S$; if $f(x) = 0$ then $g(x)$ is not zero and the same argument can be applied to $g(x)$. Let $h(x)$ be an element of minimal degree of $S$. We claim that $h(x)$ divides both $f(x)$ and $g(x)$. Since the situation is symmetric, we just prove $h(x)|f(x)$. Suppose not, then we can write $f(x) = q(x)h(x) + r(x)$, where $r(x)$ is a non-zero polynomial of degree smaller than $h(x)$. Then $r(x) = f(x) - q(x)(u(x)f(x) + v(x)g(x)) = (1 - q(x)u(x))f(x) - q(x)v(x)g(x)$ and so, if we let $r_1(x)$ be $r(x)$ divided by its leading coefficients, we see that $r_1(x) \in S$ and has degree smaller than $h(x)$, which is a contradiction.

By construction, $h(x)$ is the monic polynomial of minimal degree having such an expression. If $t(x)$ divides both $g(x)$ and $f(x)$ then $t(x)|(u(x)f(x) + v(x)g(x)) = h(x)$. Therefore, $h(x)$ is a monic polynomial of the largest possible degree dividing both $f(x), g(x)$. Suppose that $h_1(x)$ is another monic polynomial dividing $f(x)$ and $g(x)$ having the largest possible degree, i.e., the degree of $h(x)$. Then, we have $h(x) = h_1(x)b(x)$, but what we just proved. Since both polynomials have the same degree $b(x)$ must be a constant polynomial, and, then, since both are monic, $b(x) = 1$. We’ve shown the gcd is unique.

---

\[\text{l.o.t.} = \text{lower order terms.}\]
15.3. The Euclidean algorithm for polynomials.

**Theorem 15.3.** Let \( f(x), g(x) \in \mathbb{F}[x] \) be non-zero polynomials, \( g(x) = a_nx^n + \text{l.o.t.} \). If \( g(x)|f(x) \) then \( (f(x), g(x)) = a^{-1}_ng(x) \). Else, define inductively,

\[
\begin{align*}
f(x) &= q_0(x)g(x) + r_0(x), \quad \deg(r_0) < \deg(g) \\
g(x) &= q_1(x)r_0(x) + r_1(x), \quad \deg(r_1) < \deg(r_0) \\
r_0(x) &= q_2(x)r_1(x) + r_2(x), \quad \deg(r_2) < \deg(r_1) \\
&\vdots \\
r_{t-2}(x) &= q_t(x)r_{t-1}(x) + r_t(x), \quad \deg(r_t) < \deg(r_{t-1}) \\
r_{t-1}(x) &= q_{t+1}(x)r_t(x).
\end{align*}
\]

This is indeed possible, and the process always terminates. Letting \( r_1(x) = c_mx^m + \cdots + c_0, \) we have \( (f(x), g(x)) = c^{-1}_m r_1(x) \). Moreover, this algorithm also allows expressing \( (f(x), g(x)) \) in the form \( u(x)f(x) + v(x)g(x) \).

**Proof.** Each step in the process is done based on Theorem 14.1. The process must terminate because the degrees decrease.

It is easy to see that \( r_t|r_{t-1} \). Suppose we know \( r_t \) divides \( r_{t-1}, r_{t-2}, \ldots, r_a \) then, since \( r_{a-1} = q_{a+1}r_a + r_{a+1} \) we get also that \( r_t|r_{a-1} \). We conclude that \( r_t \) divides \( r_0, r_1, \ldots, r_t \). Exactly the same argument gives that \( r_t \) divides \( g(x) \) and \( f(x) \).

Conversely, if \( a(x) \) divides \( f(x) \) and \( g(x) \) then \( a(x)|(f(x)−q_0(x)g(x)) = r_0(x) \) and so \( a(x)|(g(x)−q_1(x)r_0(x)) = r_1(x) \), etc. We see that \( a(x)|r_t(x) \) and so \( r_t(x) \), once divided by its leading coefficient, must be the greatest common divisor of \( f(x) \) and \( g(x) \).

**Example 15.4.**

1. \( f(x) = x^2 + 1, g(x) = x^2 + 2ix - 1 \), complex polynomials. We have
   \[
   f(x) = 1 \cdot (x^2 + 2ix - 1) + (-2ix + 2)
   \]
   \[
   (x^2 + 2ix - 1) = \left(\frac{1}{-2i} - \frac{1}{2}\right)(-2ix + 2).
   \]
   It follows that \( (f(x), g(x)) = \frac{1}{2i}(-2ix + 2) = x + i \). This implies that \(-i\) is a root of both polynomials, as one can verify.

2. Now we choose \( \mathbb{F} = \mathbb{Z}/3\mathbb{Z} \), the field with 3 elements. We take \( f(x) = x^3 + 2x + 1, g(x) = x^2 + 1 \). We then have,
   \[
   f(x) = x \cdot (x^2 + 1) + (x + 1)
   \]
   \[
   (x^2 + 1) = (x - 1) \cdot (x + 1) + 2
   \]
   \[
   x + 1 = (2x + 2) \cdot 2.
   \]
   This implies that \( (f(x), g(x)) = 1 \). We have
   \[
   2 = (x^2 + 1) - (x - 1) \cdot (x + 1)
   \]
   \[
   = g(x) - (x - 1)(f(x) - xg(x))
   \]
   \[
   = (-x + 1)f(x) + (x^2 - x + 1)g(x).
   \]
And so we find (note that 1 = −2 in $\mathbb{F}$)
\[
1 = (f(x), g(x)) = (x - 1)f(x) - (x^2 - x + 1)g(x).
\]
(3) Consider the polynomials $f(x) = x^3 + 5x^2 + 4x, g(x) = x^3 + x^2 - x - 1$ as rational polynomials. Then
\[
f(x) = 1 \cdot g(x) + 4x^2 + 5x + 1
\]
\[
x^3 + x^2 - x - 1 = \left(\frac{1}{4}x \cdot \frac{1}{16}\right)(4x^2 + 5x + 1) - \frac{15}{16}x - \frac{15}{16}
\]
\[
(4x^2 + 5x + 1) = \frac{16}{15}(4x + 1)(\frac{15}{16}x - \frac{15}{16}).
\]
It follows that $(f(x), g(x)) = x + 1$.

To express $x + 1$ as $u(x)f(x) + v(x)g(x)$ we work backwards:
\[
\frac{-15}{16}(x + 1) = g(x) - \left(\frac{1}{4}x \cdot \frac{1}{16}\right)(4x^2 + 5x + 1)
\]
\[
= g(x) - \left(\frac{1}{4}x \cdot \frac{1}{16}\right)(f(x) - g(x))
\]
\[
= \left(-\frac{1}{4}x \cdot \frac{1}{16}\right) \cdot f(x) + \left(\frac{15}{16} + \frac{1}{4}\right)x
\]
Thus,
\[
x + 1 = (f(x), g(x)) = \left(-\frac{1}{15} + \frac{4}{15}x\right) \cdot f(x) - (1 + \frac{4}{15}x) \cdot g(x).
\]
(4) Now consider the same polynomials over the field $\mathbb{F} = \mathbb{Z}/3\mathbb{Z}$. We now have:
\[
f(x) = 1 \cdot g(x) + x^2 + 2x + 1
\]
\[
x^3 + x^2 - x - 1 = (x - 1)(x^2 + 2x + 1).
\]
Therefore, now we have $(f(x), g(x)) = x^2 + 2x + 1 = (x + 1)^2$.

15.4. Irreducible polynomials and unique factorization. Let $\mathbb{F}$ be a field. We define a relation on polynomials $f(x) \in \mathbb{F}[x]$. We say that $f(x) \sim g(x)$ if there is an element $a \in \mathbb{F}, a \neq 0$ such that $f(x) = ag(x)$.

Lemma 15.5. This relation is an equivalence relation. Related polynomials are called associates.

Proof. The relation is reflexive because $f(x) = 1 \cdot f(x)$ and symmetric, because $f(x) = ag(x)$ implies $g(x) = a^{-1}f(x)$. It is also transitive since $f(x) = ag(x)$ and $g(x) = bh(x)$ implies $f(x) = abh(x)$ and $ab \neq 0$. □

A non-constant polynomial $f$ is called irreducible if $g|f$ implies that $g \sim 1$ or $g \sim f$. That’s as close we can get to the notion of a prime. We cannot expect $g(x) = 1$, or $g(x) = f(x)$, of course. Note that if $g|h$ and $g_1 \sim g$ then $g_1|h$. 

Proposition 15.6. Let \( f(x) \in \mathbb{F}[x] \) be a non-constant polynomial. The following are equivalent:

1. \( f \) is irreducible.
2. If \( f \mid gh \) then \( f \mid g \) or \( f \mid h \).

Proof. Suppose that \( f \) is irreducible, \( f \mid gh \) and \( f \nmid g \). The only monic polynomials dividing \( f \) are 1 and \( a^{-1}f \), where \( a \) is the leading coefficient of \( f \). Therefore, \((f, g) = 1\) and so, for suitable polynomials \( u, v \) we have \( uf + vg = 1 \). Then \( uh + vgh = h \). Since \( f \) divides the left hand side, it also divides the right hand side, i.e., \( f \mid h \).

Suppose now that \( f \) has the property \( f \mid gh \Rightarrow f \mid g \) or \( f \mid h \). Let \( g \) be a divisor of \( f \). Then \( f = gh \) for some \( h \) and so \( f \mid gh \). Therefore, \( f \mid g \) or \( f \mid h \). Since \( h \mid f \), the situation is symmetric and we can assume that \( g \mid f \) and \( f \mid g \). This implies that \( \deg(g) \leq \deg(f) \) and \( \deg(f) \leq \deg(g) \), and so \( \deg(f) = \deg(g) \). But then \( \deg(h) = \deg(f) - \deg(g) = 0 \) and so \( h \) is a constant polynomial. We find that \( f \sim g \).

Example 15.7. Here are some comments on irreducible polynomials.

1. Every linear polynomial is irreducible.
2. If \( f = ax^2 + bx + c \) is reducible then \( f = (\alpha x + \beta)(\gamma x + \delta) \), where \( \alpha, \beta, \gamma, \delta \in \mathbb{F} \) and \( \alpha \neq 0, \gamma \neq 0 \). It follows then that \( f \) has a root in \( \mathbb{F} \), for example \( x = -\alpha^{-1}\beta \).

Conversely, suppose that \( f \) has a root \( \alpha \in \mathbb{F} \) then, as we shall see shortly (Theorem 15.12), \( f = (x - \alpha)g(x) \) for some polynomial \( g(x) \in \mathbb{F}[x] \), and degree considerations dictate that \( g(x) \) is a also a linear polynomial.

Therefore, for quadratic polynomials one can say that \( f \) is reducible if and only if \( f \) has a root in \( \mathbb{F} \). If, furthermore, \( 2 \neq 0 \) in the field \( \mathbb{F} \), then as we have seen in the assignment, we know that \( f \) has a root if and only if \( b^2 - 4ac \) is a square in \( \mathbb{F} \). In fact, in that case, the unique factorization of \( f \) is

\[
a x^2 + bx + c = a \left( x - \frac{-b + \sqrt{b^2 - 4ac}}{2a} \right) \left( x - \frac{-b - \sqrt{b^2 - 4ac}}{2a} \right).
\]

3. If \( f \) has degree 3 it is still true that \( f \) is reducible if and only if \( f \) has a root. But if \( f \) has degree 4 or higher this may fail. For example, the polynomial \( x^2 - 2 \) is irreducible over \( \mathbb{Q} \) because \( \sqrt{2} \) is irrational. Same for \( x^2 - 3 \). Thus, for example, the polynomials \((x^2 - 2)^2, (x^2 - 2)(x^2 - 3), (x^2 - 3)^2 \) are reducible over \( \mathbb{Q} \) but don’t have a root. (Indeed, if \( \alpha \) is a root of \((x^2 - 2)(x^2 - 3) \) then in \( \mathbb{C} \) we have \((\alpha - \sqrt{2})(\alpha + \sqrt{2})(\alpha - \sqrt{3})(\alpha + \sqrt{3}) = 0 \) and so \( \alpha = \pm \sqrt{2} \) or \( \pm \sqrt{3} \) and, in any case, is not rational.

4. The property of \( f \) being irreducible depends on the field. It is not an absolute property. For example, \( x^2 - 2 \) is irreducible in \( \mathbb{Q}[x] \) by reducible in \( \mathbb{C}[x] \) because there we can write \( x^2 = (x - \sqrt{2})(x + \sqrt{2}) \).

Theorem 15.8. (Unique factorization for polynomials) Let \( f(x) \in \mathbb{F}[x] \) be a non-zero polynomial. Then there is an \( a \in \mathbb{F} \times \) and monic irreducible polynomials \( f_1, \ldots, f_g \) of positive degree and
positive integers \( r_1, \ldots, r_g \) such that
\[
f = af_1^{r_1} \cdots f_g^{r_g}.
\]
Moreover, if
\[
f = bh_1^{s_1} \cdots h_t^{s_t},
\]
where \( b \in \mathbb{F}^\times, h_i \) monic irreducible polynomials of positive degree and \( s_i > 0 \), then \( a = b, g = t, \) and after re-naming the \( h_i \)'s we have \( h_i = f_i \) for all \( i \) and \( r_i = s_i \) for all \( i \).

**Proof.** The proof is very similar to the proof for integers. We first prove the existence of factorization. Suppose that there is a non-zero polynomial \( f(x) \) with no such factorization. Choose then a non-zero polynomial \( f(x) \) of minimal degree for which no such factorization exists. Then \( f(x) \) is not a constant polynomial and is not an irreducible polynomial either, else \( f(x) = a_n x^n + \cdots + a_0 = a_n \cdot (a_n^{-1} f(x)) \) is a suitable factorization. It follows that \( f(x) = f_1(x) f_2(x) \), where each \( f_i(x) \) has degree less than that of \( f(x) \).

Therefore, each \( f_i(x) \) has a factorization
\[
f_1(x) = c_1 a_1(x) \cdots a_m(x), \quad f_2(x) = c_2 b_1(x) \cdots b_n(x),
\]
with \( c_i \in \mathbb{F} \) and \( a_j, b_j \) monic irreducible polynomials. It follows that
\[
f(x) = (c_1 c_2) a_1(x) \cdots a_m(x) b_1(x) \cdots b_n(x),
\]
has also a factorization as claimed. Contradiction. Thus no such \( f(x) \) exists and every polynomial has a factorization as claimed.

We now show the uniqueness of the factorization. Suppose that
\[
f(x) = c_1 a_1(x) \cdots a_m(x) = c_2 b_1(x) \cdots b_n(x),
\]
with \( c_i \in \mathbb{F} \) and \( a_j, b_j \) monic irreducible polynomials. We prove the result by induction on degree \( f \).
Since \( c_i \) is the leading coefficient of \( f \), we have \( c_1 = c_2 \). In particular, the case of \( \deg(f) = 0 \) holds. Assume we have proved uniqueness for all polynomials of degree \( \leq n \) and \( \deg(f) = n + 1 \).
Since \( a_1(x) | c_2 b_1(x) \cdots b_n(x) \) and \( a_1(x) \) is irreducible, it follows that \( a_1(x) | c_2 \) (which is impossible because \( c_2 \) is a constant) or \( a_1(x) | b_i(x) \) for some \( i \) (and in particular we must have \( n \geq 1 \)). But since \( b_i(x) \) is irreducible it then follows that that \( a_1(x) \sim b_i(x) \) and so, both polynomials being monic, \( a_1(x) = b_i(x) \).

Let us re-number the \( b_i \) so that \( a_1 = b_1 \). Then, dividing by \( a_1(x) \) we have
\[
c_1 a_2(x) \cdots a_m(x) = c_2 b_2(x) \cdots b_n(x).
\]
Induction gives that \( m = n \) and, after re-numbering the \( b_i, a_i(x) = b_i(x), i = 2, 3, \ldots, n. \)

**Example 15.9.** Here are some examples.

1. \( f(x) = ax + b, a \neq 0 \) has unique factorization \( a(x + a^{-1} b) \).
(2) \( f(x) = ax^2 + bx + c \) is irreducible if and only if it has no root in \( F \), as we have seen above. If this is the case, 

\[
f(x) = a(x^2 + a^{-1}bx + a^{-1}c)
\]

is the unique factorization. Otherwise, \( f \) has two roots, say \( \beta \) and \( \gamma \) (and if \( 2 \neq 0 \) we have a formula for them) and

\[
f(x) = a(x - \alpha)(x - \beta)
\]

is the unique factorization.

(3) Consider the polynomial \( f(x) = (x^2 - 2)(x^2 - 3) \) over the field \( K = \mathbb{Q}(\sqrt{2}) \). We claim that \( x^2 - 3 \) is irreducible over \( K \). Indeed, if not, \( \sqrt{3} \in K \) and so \( \sqrt{3} = a + b\sqrt{2} \) for some rational numbers \( a, b \). Squaring, we get

\[
3 = a^2 + 2b^2 + 2ab\sqrt{2}.
\]

But this implies that \( 2ab\sqrt{2} \) is rational and so \( ab = 0 \). If \( b = 0 \) we get that \( \sqrt{3} = a \) is rational, which is a contradiction. If \( a = 0 \) we get that \( 3 = 2a^2 \), which is a contradiction, because of unique factorization of rational numbers: The power of 2 in the right hand side (i.e., in 3) is 0, while in the right hand side (i.e. in \( 2a^2 \)) is odd, whatever it may be. Therefore,

\[
f(x) = (x - \sqrt{2})(x + \sqrt{2})(x^2 - 3)
\]

is the unique factorization of \( f \) over \( \mathbb{Q}(\sqrt{2}) \).

(4) What is the unique factorization of \( x^4 + x \) over \( \mathbb{F}_2[x] \)? An obvious factor is \( x \) and then we are left with \( x^3 + 1 \). We note that \( x = 1 \) is a root (\( 2 = 0 \) now!) and so \( x^3 + 1 = (x + 1)(x^2 + x + 1) \) (cf. Theorem 15.12). The polynomial \( x^2 + x + 1 \) is quadratic and so is reducible over \( \mathbb{F}_2 \) if and only if it has a root in \( \mathbb{F}_2 \). We try and find that neither \( x = 0 \) nor \( x = 1 \) are roots. We conclude that the unique factorization is given in this case by

\[
x^4 + x = x(x + 1)(x^2 + x + 1).
\]

(5) Over the complex numbers any non-constant polynomial \( f \) factors as

\[
f(x) = a_n \prod_{i=1}^{n} (x - z_i),
\]

(see Proposition 6.7) and this is precisely its unique factorization.

We can now deduce from Theorem 15.8 the analogues of Proposition 10.6 and Corollary 10.7. The proofs are the same.

**Proposition 15.10.** Let \( f, g \) be non-zero polynomials in \( \mathbb{F}[x] \). Then \( f \mid g \) if and only if \( f = \prod_{i=1}^{m} p_i^{a_i} \) and \( g = \prod_{i=1}^{m} q_i^{b_i} \) (products of distinct irreducible monic polynomials \( p_i; a, b \) non-zero scalars) with \( a_i \geq a_i \) for all \( i = 1, \ldots, m \).
Corollary 15.11. Let $f = a_p^a_1 \cdots p_m^a_m, g = b_p^b_1 \cdots p_m^b_m$ with $p_i$ distinct irreducible monic polynomials, $a, b$ non-zero scalars and $a_i, b_i$ non-negative integers. (Any two non-zero polynomials can be written this way). Then

$$\gcd(f, g) = p_1^{\min(a_1, b_1)} \cdots p_m^{\min(a_m, b_m)}.$$ 

15.5. Roots. Let $F$ be a field and let $f(x) \in F[x]$ be a non-zero polynomial. Recall that an element $a \in F$ is called a root (or zero, or solution) of $f$ if $f(a) = 0$.

Theorem 15.12. Let $f(x) \in F[x]$ be a non-zero polynomial.

1. If $f(a) = 0$ then $f(x) = (x - a)g(x)$ for a unique polynomial $g(x) \in F[x]$. In particular, if $f$ is irreducible then $f$ has no roots in $F$.

2. Let $\deg(f) = d$ then $f$ has at most $d$ roots.

Proof. Suppose that $f(a) = 0$ and divide $f$ by $x - a$ with a residue, getting $f(x) = g(x)(x - a) + r(x)$, where $r(x)$ is either zero or a polynomial of degree less than that of $x - a$. That is, in either case, $r(x)$ is a constant. Substitute $x = a$. We get $0 = f(a) = g(a)(a - a) + r = r$ and so $f(x) = (x - a)g(x)$.

Consider the factorization of $f$ into irreducible monic polynomials:

$$f = A(x - a_1)^{s_1} \cdots (x - a_m)^{s_m}f_1(x)^{r_1} \cdots f_n(x)^{r_n},$$

where the $f_i$ are irreducible polynomials of degree larger than 1, and the $r_i, s_i$ are positive. Note that if $f(a) = 0$ then, since $f_i(a) \neq 0$ (else $f_i(x) = (x - a)g_i(x)$), we must have $a = a_i$ for some $i$. It follows that the number of roots of $f$, counting multiplicities, is $s_1 + s_2 + \cdots + s_m = \deg((x - a_1)^{s_1} \cdots (x - a_m)^{s_m}) \leq \deg(f) = d$. \hfill \Box

A field $F$ is called algebraically closed if any non-constant polynomial $f(x) \in F[x]$ has a root in $F$. The Fundamental Theorem of Algebra says that the field of complex numbers $\mathbb{C}$ is algebraically closed. It is a fact (proven in MATH370) that every field is contained in an algebraically closed field.

If $F$ is algebraically closed, then the only irreducible polynomials over $F$ are the linear ones $x - a$, $a \in F$. It follows then that

$$f(x) = A(x - a_1)^{s_1} \cdots (x - a_m)^{s_m},$$

where $A$ is the leading coefficient of $f$ and $a_1, \ldots, a_m$ are the roots (with multiplicities $s_1, \ldots, s_m$).

A natural question is, for a given field $F$ and a given polynomial $f(x)$, to tell if $f$ has a root in $F$ or not. This is in general impossible, but we have some partial answers in special cases.

Proposition 15.13. Let $f(x) = a_nx^n + \cdots + a_1x + a_0$ be a non-constant polynomial with integer coefficients. If $a = s/t$, $(s, t) = 1$, is a rational root of $f$ then $s | a_0$ and $t | a_n$. 
Proof. We have \( a_n(s/t)^n + \cdots + a_1(s/t) + a_0 = 0 \) and so
\[
a_n s^n + a_{n-1} s^{n-1} t + \cdots + a_1 s t^{n-1} + a_0 t^n = 0.
\]
Since \( s \) divides \( a_n s^n + a_{n-1} s^{n-1} t + \cdots + a_1 s t^{n-1} \), it follows that \( s | a_0 t^n \). Then, since \((s, t) = 1\), we get that \( s | a_0 \). Similarly, \( t \) divides \( a_{n-1} s^{n-1} t + \cdots + a_1 s t^{n-1} + a_0 t^n \), so \( t \) divides \( a_n s^n \). Now \((s, t) = 1\) implies that \( t | a_n \).

Example 15.14. Problem: Find the rational roots of the polynomial \( x^4 - \frac{7}{2} x^3 + \frac{5}{2} x^2 - \frac{7}{2} x + \frac{3}{2} \). The roots are the same as for the polynomial \( 2x^4 - 7x^3 + 5x^2 - 7x + 3 \). There are thus of the form \( s/t \), where \( s = \pm 1, \pm 3 \), \( t = \pm 1, \pm 2 \). We have the possibilities \( \pm 1, \pm 1/2, \pm 3, \pm 3/2 \). By checking each case, we find the roots are \( 1/2 \) and \( 3 \). We remark that after having found the root \( 1/2 \) we can divide the polynomial \( 2x^4 - 7x^3 + 5x^2 - 7x + 3 \) by \( x - 1/2 \) finding \( 2x^3 - 6x^2 + 2x - 6 \), whose roots are the roots of \( x^3 - 3x^2 + x - 3 \). So, in fact, the only possibilities for additional roots are \( \pm 3 \). We saved this way the need to check if \( \pm 3/2 \) are roots.

Here is another example. Is the polynomial \( x^3 + 2x^2 + 5 \) irreducible over \( \mathbb{Q} \)? In this case, if it is reducible then one of the factors would have to have degree 1 (this type of argument only works for degrees 1, 2, 3 polynomials. For higher degree, we might have a reducible polynomial with no linear factor, e.g., \((x^2 + 1)(x^2 + 3))\). Namely, the polynomial would have a rational root. But the rational roots can only be \( \pm 1, \pm 5 \) and one verifies those are not roots. Thus, the polynomial is irreducible.

Proposition 15.15. If \( f(x) \in \mathbb{R}[x] \) is a polynomial of odd degree then \( f \) has a root in \( \mathbb{R} \).

Proof. Since the roots of \( f \) are the roots of \( -f \), we may assume that \( f(x) = a_n x^n + \cdots + a_1 x + a_0 \), \( a_i \in \mathbb{R}, a_n > 0 \). An easy estimate shows that there is an \( N > 0 \) such that \( f(N) > 0 \) and \( f(-N) < 0 \). By the intermediate value theorem there is some \( a, -N < a < N \) such that \( f(a) = 0 \). \( \square \)

15.6. Eisenstein’s criterion. (Planned)

15.7. Roots of polynomials in \( \mathbb{Z}/p\mathbb{Z} \). Let \( p \) be a prime and let \( \mathbb{Z}/p\mathbb{Z} \) be the field with \( p \) elements whose elements are congruence classes modulo \( p \). By Fermat’s little theorem, every element of \( \mathbb{Z}/p\mathbb{Z}^\times \) is a root of \( x^{p-1} - 1 \). This gives \( p - 1 \) distinct roots of \( x^{p-1} - 1 \) and so these must be all the roots and each with multiplicity one. It follows that the roots of \( x^p - x \) are precisely the elements of \( \mathbb{Z}/p\mathbb{Z} \) again with multiplicity one. That is,
\[
x^p - x = \prod_{a=0}^{p-1} (x - a).
\]

Proposition 15.16. Let \( f(x) \) be any polynomial in \( \mathbb{Z}/p\mathbb{Z}[x] \). Then \( f(x) \) has a root in \( \mathbb{Z}/p\mathbb{Z} \) if and only if \( \gcd(f(x), x^p - x) \neq 1 \).
Proof. If $f(a) = 0$ for some $a \in \mathbb{Z}/p\mathbb{Z}$ then $(x - a)|f(x)$, but also $(x - a)|(x^p - x)$. It follows that $\gcd(f(x), x^p - x) \neq 1$. Conversely, if $h(x) = \gcd(f(x), x^p - x) \neq 1$ then, since $h(x)|x^p - x = \prod_{a=0}^{p-1}(x - \bar{a})$, by unique factorization, we must have $h(x) = \prod_{i=1}^{n}(x - a_i)$ for some distinct elements $a_1, \ldots, a_n$ of $\mathbb{Z}/p\mathbb{Z}$. In particular, each such $a_i$ is a root of $f(x)$. \qed

The straightforward way to check if $f(x)$ has a root in $\mathbb{Z}/p\mathbb{Z}$ is just to try all possibilities for $x$. Suppose that $f(x)$ has a small degree relative to $p$, say 5, to fix ideas. We still have to try $p$ numbers, each in its turn, to see if any of which is a root. But $p$ may be very large, much too large for this method to be feasible. For example, $p$ might be of cryptographic size $\approx 2^{512}$. Even with a computer doing $10^{10}$ operations per second, checking all these possibilities will take more than $10^{134}$ years!

Proposition 15.16 suggests a different method: Calculate $\gcd(f(x), x^p - x)$. Note that except for the first step

$$x^p - x = q_0(x)f(x) + r_0(x),$$

all the polynomials involved in the Euclidean algorithm would have very small degrees (smaller then $f$’s for example) and so the Euclidean algorithm will terminate very quickly. The first step, though, could be very time consuming given what we know at this point. Later we shall see that it can, in fact, be done quickly.

We have seen that many of the features of arithmetic in $\mathbb{Z}$ can be carried out in $\mathbb{F}[x]$. We still don’t have an analogue of passing from $\mathbb{Z}$ to $\mathbb{Z}/n\mathbb{Z}$ in the context of $\mathbb{F}[x]$. This is one of our main motivation for studying rings in much more detail.
Part 5. Rings

16. SOME BASIC DEFINITIONS AND EXAMPLES

Recall our definition of a ring

**Definition 16.1.** A ring $R$ is a non-empty set together with two operations, called “addition” and “multiplication” that are denoted, respectively, by

$$(x, y) \mapsto x + y, \quad (x, y) \mapsto xy.$$ 

One requires the following axioms to hold:

1. $x + y = y + x, \forall x, y \in R$. (Commutativity of addition)
2. $(x + y) + z = x + (y + z), \forall x, y, z \in R$. (Associativity of addition)
3. There exists an element in $R$, denoted 0, such that $0 + x = x, \forall x \in R$. (Neutral element for addition)
4. $\forall x \in R, \exists y \in R$ such that $x + y = 0$. (Inverse with respect to addition)
5. $(xy)z = x(yz), \forall x, y, z \in R$. (Associativity of multiplication)
6. There exists an element 1 $\in R$ such that $1x = x, \forall x \in R$. (Neutral element for multiplication)
7. $z(x + y) = zx + zy, (x + y)z = xz + yz, \forall x, y, z \in R$. (Distributivity)

Recall also that a ring $R$ is called a division ring (or sometimes a skew-field) if $1 \neq 0$ in $R$ and any non-zero element of $R$ has an inverse w.r.t. multiplication. A commutative division ring is precisely what we call a field.

**Example 16.2.** $\mathbb{Z}$ is a commutative ring. It is not a division ring and so is not a field. The rational numbers $\mathbb{Q}$ form a field. The real numbers $\mathbb{R}$ form a field. The complex numbers $\mathbb{C}$ form a field.

We have also noted some useful formal consequences of the axioms defining a ring:

1. The element 0 appearing in axiom (3) is unique.
2. The element $y$ appearing in axiom (4) is unique. We shall denote $y$ by $-x$.
3. We have $-(-x) = x$ and $-(x + x') = -x - x'$, where, technically $-x - x'$ means $(-x) + (-x')$.
4. We have $x \cdot 0 = 0, 0 \cdot x = 0$.

Here are some further examples. We do not prove the ring axioms hold; this is left as an exercise.

**Example 16.3.** Let $\mathbb{F}$ be a field and $n \geq 1$ an integer. Consider the set of $n \times n$ matrices:

$$M_n(\mathbb{F}) = \left\{ \begin{pmatrix} a_{11} & \cdots & a_{1n} \\ \vdots & \ddots & \vdots \\ a_{n1} & \cdots & a_{nn} \end{pmatrix} : a_{ij} \in \mathbb{F} \right\}.$$ 

For example:
(1) for \( n = 1 \) we just get \((a_{11}), a_{11} \in \mathbb{F}\);

(2) for \( n = 2 \), \[
\begin{pmatrix}
a_{11} & a_{12} \\
a_{21} & a_{22}
\end{pmatrix};
\]

(3) for \( n = 3 \) we get \[
\begin{pmatrix}
a_{11} & a_{12} & a_{13} \\
a_{21} & a_{22} & a_{23} \\
a_{31} & a_{32} & a_{33}
\end{pmatrix}.
\]

In general we shall write an \( n \times n \) matrix as \((a_{ij})\), or \((a_{ij})_{i,j=1}^n\). The index \( i \) is the row index and the index \( j \) is the column index. We then define

\[
(a_{ij}) + (b_{ij}) = (a_{ij} + b_{ij}), \quad (a_{ij})(b_{ij}) = (c_{ij}),
\]

where

\[
c_{ij} = \sum_{k=1}^n a_{ik} b_{kj}.
\]

For example:

(1) for \( n = 1 \) we get \((a) + (b) = (a + b) \) and \((a)(b) = (ab)\). Namely, we just get \( \mathbb{F} \) again!

(2) for \( n = 2 \), we have

\[
\begin{pmatrix}
a_{11} & a_{12} \\
a_{21} & a_{22}
\end{pmatrix} + \begin{pmatrix}
b_{11} & b_{12} \\
a_{21} & b_{22}
\end{pmatrix} = \begin{pmatrix}
a_{11} + b_{11} & a_{12} + b_{12} \\
a_{21} + b_{21} & a_{22} + b_{22}
\end{pmatrix},
\]

and

\[
\begin{pmatrix}
a_{11} & a_{12} \\
a_{21} & a_{22}
\end{pmatrix} \begin{pmatrix}
b_{11} & b_{12} \\
b_{21} & b_{22}
\end{pmatrix} = \begin{pmatrix}
a_{11}b_{11} + a_{12}b_{21} & a_{11}b_{12} + a_{12}b_{22} \\
a_{21}b_{11} + a_{22}b_{21} & a_{21}b_{12} + a_{22}b_{22}
\end{pmatrix}.
\]

Under these definitions \( M_n(\mathbb{F}) \) is a ring, called the ring of \( n \times n \) matrices with entries in \( \mathbb{F} \), with identity given by the identity matrix

\[
I_n = \begin{pmatrix}
1 & 0 & \ldots & 0 \\
0 & 1 & \ldots & 0 \\
0 & 0 & \ldots & 1
\end{pmatrix},
\]

and zero given by the zero matrix (the matrix all whose entries are zero). For \( n \geq 2 \) this is a non-commutative ring. Indeed, for \( n = 2 \), for example, we have,

\[
\begin{pmatrix}
1 & 1 \\
0 & 1
\end{pmatrix} \begin{pmatrix}
1 & 0 \\
1 & 1
\end{pmatrix} = \begin{pmatrix}
2 & 1 \\
1 & 1
\end{pmatrix}, \quad \begin{pmatrix}
1 & 0 \\
1 & 1
\end{pmatrix} \begin{pmatrix}
1 & 1 \\
0 & 1
\end{pmatrix} = \begin{pmatrix}
1 & 1 \\
1 & 2
\end{pmatrix}.
\]

These are never equal, else \( 2 = 1 \) in \( \mathbb{F} \), which implies \( 1 = 0 \) in \( \mathbb{F} \), which is never the case, by definition.

**Example 16.4.** Let \( \epsilon \) be a formal symbol and \( \mathbb{F} \) a field. The ring of dual numbers, \( \mathbb{F}[\epsilon] \), is defined as

\[
\mathbb{F}[\epsilon] = \{ a + b\epsilon : a, b \in \mathbb{F} \},
\]
with the following addition and multiplication:

\[(a + bc) + (c + de) = a + c + (b + d)e, \quad (a + be)(c + de) = ac + (ad + bc)e.\]

Note that \(e\) is a zero divisor: \(e \neq 0\) but \(e^2 = 0\).

**Example 16.5.** Let \(n \neq \pm 1\) be a square free integer (i.e., if \(p|n\), \(p\) prime, then \(p^2 \nmid n\)). Then \(\sqrt{n}\) is not a rational number. Indeed if \(\sqrt{n}\) is rational, \(\sqrt{n} = s/t\), \((s, t) = 1\), then \(n = s^2/t^2\). Let \(p\) a prime dividing \(n\) and so \(p^2 \nmid n\). Then \(nt^2 = s^2\) and so \(p|s^2\). But then \(p|s\). Looking at the power of \(p\) in the unique factorization of both sides, it follows that \(p|t\) and thus \(p|(s, t) - a\) contradiction.

Consider

\[\mathbb{Z} [\sqrt{n}] = \{a + b\sqrt{n} : a, b \in \mathbb{Z}\}.\]

This is a subset of \(\mathbb{C}\), containing 0 and 1 and is closed under addition and multiplication:

\[(a + b\sqrt{n}) + (c + d\sqrt{n}) = a + c + (b + d)\sqrt{n}, \quad (a + b\sqrt{n})(c + d\sqrt{n}) = ac + bdn + (ad + bc)\sqrt{n}.\]

We remark that any element of this ring has a unique expression as \(a + b\sqrt{n}\). Indeed, if \(a + b\sqrt{n} = c + d\sqrt{n}\), either \(b = d\) (and then obviously \(a = c\)) or \(\sqrt{n} = (a - c)/(d - b)\) is a rational number, which it’s not.

**Example 16.6.** Let \(R_1, R_2\) be rings. Then \(R_1 \times R_2\) is a ring with the following operations:

\[(a_1, b_1) + (a_2, b_2) = (a_1 + a_2, b_1 + b_2), \quad (a_1, b_1)(a_2, b_2) = (a_1a_2, b_1b_2).\]

The zero element is \((0_{R_1}, 0_{R_2})\) and the identity element is \((1_{R_1}, 1_{R_2})\). The ring \(R_1 \times R_2\) is called the **direct product** of \(R_1\) and \(R_2\).

**Definition 16.7.** Let \(S \subset R\) be a subset. \(S\) is called a **subring** of \(R\) if the following holds:

1. \(0_R, 1_R\) belong to \(S\);
2. \(s_1, s_2 \in S \Rightarrow s_1 \pm s_2 \in S\);
3. \(s_1, s_2 \in S \Rightarrow s_1s_2 \in S\).

Note that in this case \(S\) is a ring in its own right.

**Example 16.8.** The easiest examples are \(\mathbb{Z}, \mathbb{Q}, \mathbb{R}\) being subrings of \(\mathbb{C}\). We’ve already seen examples of subrings of the ring of \(2 \times 2\) matrices in the assignments.

Consider the subset \(\{(r, 0) : r \in \mathbb{R}\}\) of the ring \(\mathbb{R} \times \mathbb{R}\). It is closed under addition and multiplication. It is even a ring because \((r, 0)(1, 0) = (r, 0)\) and so \((1, 0)\) serves as an identity element for this subset. Nonetheless, it is not a subring of \(\mathbb{R} \times \mathbb{R}\), because the identity element of \(\mathbb{R} \times \mathbb{R}\), which is \((1, 1)\), does not belong to the set \(\{(r, 0) : r \in \mathbb{R}\}\).
17. Ideals

Definition 17.1. Let $R$ be a ring. A (two-sided) ideal $I$ of $R$ is a subset of $R$ such that

1. $0 \in I$;
2. if $a, b \in I$ then $a + b \in I$;
3. if $a \in I$, $r \in R$, then $ra \in I$ and $ar \in I$.

Remark 17.2. Note that if $a \in I$ then $-a = -1 \cdot a \in I$.

We shall use the notation $I \triangleleft R$ to indicate that $I$ is an ideal of $R$.

Example 17.3. $I = \{0\}$ and $I = R$ are always ideals. They are called the trivial ideals.

Example 17.4. Suppose that $R$ is a division ring (e.g., a field) and $I \triangleleft R$ is a non-zero ideal. Then $I = R$. Indeed, there is an element $a \in I$ such that $a \neq 0$. Then $1 = a^{-1}a \in I$ and so for every $r \in R$ we have $r = r \cdot 1 \in I$. That is, $I = R$. We conclude that a division ring has only the trivial ideals. (Note also that the argument shows for any ring $R$ that if an ideal $I$ contains an invertible element of $R$ then $I = R$.)

Example 17.5. Let $R$ be a commutative ring. Let $r \in R$. The principal ideal $(r)$ is defined as

$$(r) = \{ra : a \in R\} = \{ar : a \in R\}.$$ 

We also denote this ideal by $rR$ or $Rr$. This is indeed an ideal: First $0 = r \cdot 0$ is in $(r)$. Second, given two elements $ra_1, ra_2$ in $(r)$ we have $ra_1 + ra_2 = r(a_1 + a_2) \in (r)$ and for every $s \in R$ we have $s(ra_1) = (sr)a_1 = (rs)a_1 = r(sa_1) \in R$ (using commutativity!), $(ra_1)s = r(a_1s) \in R$.

Definition 17.6. Let $R$ be a commutative ring. If every ideal of $R$ is principal, one calls $R$ a principal ideal ring.

Theorem 17.7. $\mathbb{Z}$ is a principal ideal ring. In fact, the list

$$(0), (1), (2), (3), (4), \ldots$$

is a complete list of the ideals of $\mathbb{Z}$. (Note that another notation is $0, 2\mathbb{Z}, 3\mathbb{Z}, 4\mathbb{Z}, \ldots$)

Proof. We already know these are ideals and we note that for $i > 0$ the minimal positive number in the ideal $(i)$ is $i$. Thus, these ideals are distinct.

Let $I$ be an ideal of $\mathbb{Z}$. If $I = \{0\}$ then $I$ appears in the list above. Else, there is some non-zero element $a \in I$. If $a < 0$ then $-a = -1 \cdot a \in I$ and so $I$ has a positive element in it. Choose the smallest positive element in $I$ and call it $i$.

First, since $i \in I$ so is $ia$ for any $a \in \mathbb{Z}$ and so $(i) \subseteq I$. Let $b \in I$. Divide $b$ by $i$ with residue: $b = qi + r$, where $0 \leq r < i$. Note that $r = b - qi$ is an element of $I$, smaller than $i$. The only possibility is that $r = 0$ and so $b \in (i)$. Thus, $I = (i)$.

Theorem 17.8. Let $\mathbb{F}$ be a field. The ring $\mathbb{F}[x]$ is a principal ideal ring. Two ideals $(f(x)), (g(x))$ are equal if and only if $f \sim g$. 


Proof. The proof is very similar to the case of \( \mathbb{Z} \). Let \( I \) be an ideal. If \( I = \{0\} \) then \( I = (0) \), the principal ideal generated by 0. Else, let \( f(x) \in I \) be a non-zero polynomial whose degree is minimal among all non-zero elements of \( I \). On the one hand \( I \supseteq \langle f(x) \rangle \). On the other hand, let \( g(x) \in I \) and write \( g(x) = q(x)f(x) + r(x) \), where \( r(x) \) is either zero or of degree small than \( f \)'s. But \( r(x) = g(x) - q(x)f(x) \in I \). Thus, we must have \( r(x) = 0 \) and so \( g(x) = q(x)f(x) \in \langle f(x) \rangle \). That is, \( I \subseteq \langle f(x) \rangle \).

Suppose that \( \langle f(x) \rangle \supseteq \langle g(x) \rangle \) then \( g(x) = f(x)h(x) \) for some polynomial \( h(x) \in \mathbb{F}[x] \). That is \( f(x)|g(x) \). Thus, if \( \langle f(x) \rangle = \langle g(x) \rangle \) then \( f|g \) and \( g|f \) and so \( f \sim g \).

If \( f|g \), say \( g(x) = f(x)h(x) \) then any multiple of \( g(x) \), say \( g(x)t(x) \) is equal to \( f(x)[h(x)t(x)] \) and so \( \langle g(x) \rangle \subseteq \langle f(x) \rangle \). If \( f \sim g \) then \( f|g \) and \( g|f \) and so, by the argument above, \( \langle f(x) \rangle = \langle g(x) \rangle \). \( \square \)

**Example 17.9.** Let \( \mathbb{F} \) be a field. One can show that all the ideals of \( \mathbb{F}[\epsilon] \) are \( \{0\} = (0), \mathbb{F}[\epsilon] = (1) \) and \( (\epsilon) = \{be : b \in \mathbb{F}\} \) and so the ring of dual numbers is also a principal ideal ring.

**Example 17.10.** The ring of polynomials \( \mathbb{C}[x, y] \) in two variables with complex coefficients is not a principal ideal ring. We claim that the set of polynomials \( I = \{f(x, y) : f(0, 0) = 0\} \), namely, polynomials with zero constant term, is an ideal that is not principal. We leave that as an exercise.

**Example 17.11.** Let \( R_1, R_2 \) be rings with ideals \( I_1, I_2 \), respectively. Then \( I_1 \times I_2 \) is an ideal of \( R_1 \times R_2 \).

**Example 17.12.** Let us consider the ring \( \mathbb{F}[x] \) and in it the set
\[
S = \{f(x) : f(x) = a_0 + a_2x^2 + a_3x^3 + \ldots\},
\]
of polynomials with no \( x \) term. Note that \( 0 \in S \) and \( s_1, s_2 \in S \Rightarrow s_1 + s_2 \in S \). However, \( S \) is not an ideal. We have \( 1 \in S \) but \( x \cdot 1 \notin S \).

**Example 17.13.** Consider the ring \( \mathbb{Z}[\sqrt{5}] \). In this ring we consider
\[
I = \{5a + b\sqrt{5} : a, b \in \mathbb{Z}\}.
\]
We claim that \( I \) is an ideal. Clearly \( 0 \in I \) and \( I \) is closed under addition. To show \( I \) is an ideal it is enough to show that for every \( x \in I \) and \( c \in \mathbb{Z} \) also \( cx \) and \( \sqrt{5}x \) are in \( I \). Indeed then for any element of the ring \( c + d\sqrt{5} \) we have \( (c + d\sqrt{5})x = cx + d(\sqrt{5}x) \); \( \sqrt{5}x \in I \) so \( d\sqrt{5}x \in I \) and also \( cx \in I \) and so \( cx + d\sqrt{5}x \in I \).

The verification is easy: \( c(5a + b\sqrt{5}) = 5(ac) + (bc)\sqrt{5} \) and so is in \( I \). Also \( \sqrt{5}(5a + b\sqrt{5}) = 5b + (5a)\sqrt{5} \) is in \( I \).

**Example 17.14.** Let \( R \) be a ring and \( I_1, I_2 \) two ideals of \( R \). Then
\[
I_1 + I_2 = \{i_1 + i_2 : i_1 \in I_1, i_2 \in I_2\}
\]
is an ideal of \( R \). Inductively, the sum of \( n \) ideals \( I_1 + I_2 + \cdots + I_n \) is an ideal. A particular case is the following: Let \( R \) be a commutative ring and \( I_i = r_iR \) a principal ideal. Then
\[
r_1R + r_2R + \cdots + r_nR
\]
is an ideal of \( R \); we often denote it by \((r_1, r_2, \ldots, r_n)\) or \(\langle r_1, r_2, \ldots, r_n \rangle\) (so in particular a principle ideal \((r)\) may also be denoted \(\langle r \rangle\)).

Let us consider the situation of the ring \( R = \mathbb{Z}[\sqrt{5}] \) and the ideal \(\langle 2, 1 + \sqrt{5} \rangle\). We know abstractly that this is an ideal. I claim that this ideal is not principle. In particular, this shows that this ideal is not \(\mathbb{Z}[\sqrt{5}]\) and, more importantly, gives us an example of a ring with non-principal ideals.

Suppose that \(\langle 2, 1 + \sqrt{5} \rangle = \langle a + b\sqrt{5} \rangle\). It follows that \(2 = (a + b\sqrt{5})(c + d\sqrt{5})\) and so that \(2 = (a - b\sqrt{5})(c - d\sqrt{5})\) (check!). Therefore, by multiplying these two equations, \(4 = (a^2 + 5b^2)(c^2 + 5d^2)\). This is an equation in integers and so \(a \in \{\pm 1, \pm 2\}, b = 0\) and we conclude that \(\langle 2, 1 + \sqrt{5} \rangle = \langle a \rangle\) is equal to \(\langle 1 \rangle\) or \(\langle 2 \rangle\). Now, if \(\langle 2, 1 + \sqrt{5} \rangle = \langle 2 \rangle\) this implies that \(1 + \sqrt{5} = 2(c + d\sqrt{5})\), which is a contradiction. If \(\langle 2, 1 + \sqrt{5} \rangle = \langle 1 \rangle\) then \(1 = 2(c_1 + d_1\sqrt{5}) + (1 + \sqrt{5})(c_2 + d_2\sqrt{5}) = (2c_1 + c_2 - 5d_2 + \sqrt{5}(2d_1 + c_2 + d_2)\). Therefore, \(2d_1 + c_2 + d_2 = 0\), that is, \(-2d_1 - c_2 = d_2\) and we get \(1 = 2c_1 + c_2 - 5d_2 = 2c_1 + c_2 + 10d_1 + 5c_2 = 2(c_1 + 3c_2 + 5d_1)\). This is an equation in integers and it implies that \(1\) is even. Contradiction.

18. Homomorphisms

Let \( R, S \) be rings. A function \( f : R \to S \) is a ring homomorphism if the following holds:

1. \( f(1_R) = 1_S \);
2. \( f(r_1 + r_2) = f(r_1) + f(r_2) \);
3. \( f(r_1 r_2) = f(r_1) f(r_2) \).

Here are some formal consequences (that are nonetheless very useful).

- \( f(0_R) = 0_S \). Indeed, \( f(0_R) = f(0_R + 0_R) = f(0_R) + f(0_R) \). Let \( y = f(0_R) \) then \( y = y + y \). Adding \(-y\) to both sides we find \( 0_S = f(0_R) \).
- We have \( f(-r) = -f(r) \). Indeed: \( 0_S = f(0_R) = f(r + (-r)) = f(r) + f(-r) \) and so \( f(-r) = -f(r) \) (just because it sums with \( f(r) \) to \( 0_S \)!)\)
- We have \( f(r_1 - r_2) = f(r_1) - f(r_2) \), because \( f(r_1 - r_2) = f(r_1 + (-r_2)) \) (this, by definition) and so \( f(r_1 - r_2) = f(r_1) + f(-r_2) = f(r_1) - f(r_2) \).

Note, in particular, that \( f(0_R) = 0_S \) is a consequence of axioms (2), (3). On the other hand \( f(1_R) = 1_S \) does not follow from (2), (3) and we therefore include it as an axiom (though not all authors do that). Here is an example. Consider,

\[
 f : \mathbb{R} \to \mathbb{R} \times \mathbb{R}, \quad f(r) = (r, 0).
\]

This map satisfies \( f(r_1 + r_2) = f(r_1) + f(r_2) \) and \( f(r_1 r_2) = f(r_1) f(r_2) \), but \( f(1) = (1, 0) \) is not the identity element of \( \mathbb{R} \times \mathbb{R} \). So this is not a ring homomorphism.

On the other hand, if \( S \subset R \) is a subring then the inclusion map \( i : S \to R, i(s) = s \), is a ring homomorphism. Note that this explains why in the definition of a subring we insisted on \( 1_R \in S \).
Proposition 18.1. Let $f : R \to S$ be a homomorphism of rings. The image of $f$ is a subring of $S$.

Proof. As we saw, $f(0_R) = 0_S$. Also, by definition $f(1_R) = 1_S$ and so $0_S, 1_S \in \text{Im}(f)$. Let now $s_1, s_2 \in \text{Im}(f)$, say $s_i = f(r_i)$. Then, $s_1 \pm s_2 = f(r_1) \pm f(r_2) = f(r_1 \pm r_2)$ and so $s_1 \pm s_2 \in \text{Im}(f)$. Similarly, $s_1 s_2 = f(r_1 r_2)$ and so $s_1 s_2 \in \text{Im}(f)$. \hfill \square

Definition 18.2. Let $f : R \to S$ be a homomorphism of rings. The kernel of $f$, $\text{Ker}(f)$, is defined as follows:

$$\text{Ker}(f) = \{ r \in R : f(r) = 0 \}.$$

Proposition 18.3. $\text{Ker}(f)$ is an ideal of $R$. The map $f$ is injective if and only if $\text{Ker}(f) = \{0\}$.

Proof. First, since $f(0_R) = 0_S$ we have $0_R \in \text{Ker}(f)$. Suppose that $r_1, r_2 \in \text{Ker}(f)$ then $f(r_i) = 0_S$ and we find that $f(r_1 + r_2) = f(r_1) + f(r_2) = 0_S + 0_S = 0_S$ so $r_1 + r_2 \in \text{Ker}(f)$.

Now suppose that $r_1 \in \text{Ker}(f)$ and $r \in R$ is any element. We need to show that $rr_1, r_1 r \in \text{Ker}(f)$. We calculate $f(rr_1) = f(r)f(r_1) = f(r)0_S = 0_S$, so $rr_1 \in \text{Ker}(f)$. Similarly for $r_1 r$.

We have so far proven that $\text{Ker}(f)$ is an ideal. Suppose now that $f$ is injective. Then $f(r) = 0_S$ implies $f(r) = f(0_R)$ and so $r = 0_R$. That is, $\text{Ker}(f) = \{0_R\}$.

Suppose conversely that $\text{Ker}(f) = \{0_R\}$. If $f(r_1) = f(r_2)$ then $0_S = f(r_1) - f(r_2) = f(r_1 - r_2)$ and so $r_1 - r_2 \in \text{Ker}(f)$. Since $\text{Ker}(f) = \{0_R\}$, we must have $r_1 - r_2 = 0_R$; that is, $r_1 = r_2$. We proved that $f$ is injective. \hfill \square

We now look at some examples:

Example 18.4. Let $n \geq 1$ be an integer. Define a function,

$$f : \mathbb{Z} \to \mathbb{Z}/n\mathbb{Z},$$

by $f(a) = \overline{a}$ (the congruence class of $a$ modulo $n$). Then $f$ is a homomorphism:

1. $f(1) = \overline{1}$ and $\overline{1}$ is the indeed the identity element of $\mathbb{Z}/n\mathbb{Z}$;
2. $f(a + b) = \overline{a + b} = \overline{a} + \overline{b} = f(a) + f(b)$;
3. $f(ab) = \overline{ab} = \overline{a} \overline{b} = f(a)f(b)$.

The kernel of $f$ is $\{ a : \overline{a} \equiv 0 \pmod{n} \} = \{ a : n|a \} = \{n\}$.

Example 18.5. Let $R, S$ be any rings and define

$$f : R \times S \to R, \quad f((r, s)) = r.$$

This is a homomorphism:

1. $f(1_{R \times S}) = f((1_R, 1_S)) = 1_R$;
2. $f((r_1, s_1) + (r_2, s_2)) = f((r_1 + r_2, s_1 + s_2)) = r_1 + r_2 = f((r_1, s_1)) + f((r_2, s_2))$;
3. $f((r_1, s_1)(r_2, s_2)) = f((r_1 r_2, s_1 s_2)) = r_1 r_2 = f((r_1, s_1))f((r_2, s_2))$.

The kernel of $f$ is $\{ (r, s) : r = 0 \} = \{ (0, s) : s \in S \} = \{0\} \times S$. 
Example 18.6. Let $\mathbb{F}$ be a field and $\mathbb{F}[\epsilon]$ the ring of dual numbers. Define

$$f : \mathbb{F}[\epsilon] \to \mathbb{F}, \quad f(a + b\epsilon) = a.$$ 

Then $f$ is a homomorphism:

1. $f(1) = 1$;
2. $f((a + b\epsilon) + (c + d\epsilon)) = f(a + c + (b + d)\epsilon) = a + c = f(a + b\epsilon) + f(c + d\epsilon)$;
3. $f((a + b\epsilon)(c + d\epsilon)) = f(ac + (ad + bc)\epsilon) = ac = f(a + b\epsilon)f(c + d\epsilon)$.

The kernel of $f$ is $\{a + b\epsilon : a = 0, b \in \mathbb{F}\} = \{be : b \in \mathbb{F}\}$. We claim that this is the ideal $(\epsilon)$. On the one hand $be$ certainly is in $(\epsilon)$ for any $b$. That is $\text{Ker}(f) \subseteq (\epsilon)$. On the other hand $(c + d\epsilon)\epsilon = c\epsilon$ and that shows $(\epsilon) \subseteq \text{Ker}(f)$.

Example 18.7. Let $\mathbb{F}$ be a field. Let $a \in \mathbb{F}$ be a fixed element. Define

$$\alpha : \mathbb{F}[x] \to \mathbb{F}, \quad \alpha(g(x)) = g(a).$$

Then $\alpha$ is a homomorphism:

1. $\alpha(1)$ is the value of the constant polynomial 1 at $a$ which is just 1, so $\alpha(1) = 1$.
2. We have $\alpha(f + g) = (f + g)(a) = f(a) + g(a) = \alpha(f) + \alpha(g)$;
3. Similarly, $\alpha(fg) = (fg)(a) = f(a)g(a) = \alpha(f)\alpha(g)$.

Therefore, $\alpha$ is a homomorphism. It is called the specialization homomorphism or the evaluation homomorphism. The kernel of $\alpha$ is $\{f \in \mathbb{F}[x] : f(a) = 0\}$ and is equal to the principal ideal $(x-a)$. Indeed: if $g(x) \in (x-a)$ then $g(x) = (x-a)g_1(x)$ and so $g(a) = (a-a)g_1(a) = 0$. Conversely, if $g(a) = 0$, Theorem 15.12 says that $g(x) = (x-a)g_1(x)$ for some polynomial $g_1(x)$ and so $g(x) \in (x-a)$.

Example 18.8. Let $A$ be the set of all continuous functions $f : [0, 1] \to \mathbb{R}$. Define the sum (resp. product) of two functions $f, g$ be the function $f + g$ (resp. $fg$) whose value at any $x$ is $f(x) + g(x)$ (resp. $f(x)g(x)$). That is:

$$(f + g)(x) = f(x) + g(x), \quad (fg)(x) = f(x)g(x).$$

This is a ring (in particular, these are operations – the sum and product of continuous functions is continuous!). Its zero element is the constant function zero and its identity element is the constant function 1. Let $x \in [0, 1]$ be a fixed element. Define

$$\varphi : A \to \mathbb{R}, \quad \varphi(f) = f(x).$$

Then $\varphi$ is a ring homomorphism whose kernel are all the functions vanishing at the point $x$. 
18.1. Units. Let $R$ be any ring. The units of $R$ are denoted $R^\times$ and defined as follows:

$$R^\times = \{x \in R : \exists y \in R, xy = yx = 1\}.$$

For example, $1_R$ is always a unit. If $R$ is a field then, by definition, $R^\times = R - \{0\}$.

**Lemma 18.9.** We have the following properties:

1. If $r_1, r_2 \in R^\times$ then $r_1 r_2 \in R^\times$.
2. Let $f : R \to S$ be a homomorphism of rings then $f(R^\times) \subseteq S^\times$.

**Proof.** Suppose that $r_1, r_2 \in R^\times$ and $r_1 y_1 = y_1 r_1 = 1, r_2 y_2 = y_2 r_2 = 1$. Let $y = y_2 y_1$ then $(r_1 r_2)y = r_1 (r_2 y_2) y_1 = r_1 \cdot 1 \cdot y_1 = r_1 y_1 = 1$. A similar computation gives $y(r_1 r_2) = 1$ and so $r_1 r_2 \in R^\times$.

Let now $f : R \to S$ be a homomorphism and $r \in R^\times$ with $ry = yr = 1_R$. Then $f(r) f(y) = f(ry) = f(1_R) = 1_S$ and $f(y) f(r) = f(yr) = f(1_R) = 1_S$. It follows that $f(r) \in S^\times$. \(\square\)

**Example 18.10.** We have $\mathbb{Z}^\times = \{\pm 1\}$. We have $\mathbb{Q}^\times = \mathbb{Q} - \{0\}$.

**Example 18.11.** We have $\mathbb{F}[\epsilon] = \{a + b \epsilon : a \neq 0\}$. Indeed, if $a \neq 0$ then $(a + b \epsilon)(a^{-1} - a^{-2} b \epsilon) = 1$ (where $a^{-2}$ is by definition $(a^2)^{-1}$. It satisfies $a^{-2} a = a^{-1}$). Conversely, if $(a + b \epsilon)(c + d \epsilon) = 1$ then $ac = 1$ and so $a \neq 0$.

**Example 18.12.** Let $n \neq \pm 1$ be a square free integer. We have

$$\mathbb{Z}[\sqrt{n}]^\times = \{a + b \sqrt{n} : a^2 - b^2 n = \pm 1\}.$$

Indeed, if $a^2 - b^2 n = \pm 1$ then $(a + b \sqrt{n})(a - b \sqrt{n}) = \pm 1$ and so $a + b \sqrt{n}$ is invertible with inverse $\pm (a - b \sqrt{n})$. Conversely, if $a + b \sqrt{n}$ is invertible, say $(a + b \sqrt{n})(c + d \sqrt{n}) = 1$ (for some $c, d \in \mathbb{Z}$) then $ad + bc = 0$ and so also $(a - b \sqrt{n})(c - d \sqrt{n}) = 1$. We get that

$$(a + b \sqrt{n})(a - b \sqrt{n})(c + d \sqrt{n})(c - d \sqrt{n}) = 1.$$

But $(a + b \sqrt{n})(a - b \sqrt{n}) = a^2 - b^2 n$ and $(c + d \sqrt{n})(c - d \sqrt{n}) = c^2 - d^2 n$ are integers. So

$$(a + b \sqrt{n})(a - b \sqrt{n}) = a^2 - b^2 n = \pm 1.$$

**Example 18.13.** Let $\mathbb{F}$ be a field. The units of the ring $M_2(\mathbb{F})$ are the matrices

$$\left\{ \begin{pmatrix} a & b \\ c & d \end{pmatrix} : ad - bc \neq 0 \right\}.$$

Indeed, suppose that for the matrix $\begin{pmatrix} a & b \\ c & d \end{pmatrix}$ we have $ad - bc \neq 0$. Consider the matrix

$$(ad - bc)^{-1} \begin{pmatrix} d & -b \\ -c & a \end{pmatrix}$$
(where by \(t\begin{pmatrix}a & b \\c & d \end{pmatrix}\) we mean \(\begin{pmatrix}ta & tb \\tc & td \end{pmatrix}\). It is equal to \(\begin{pmatrix}a & b \\c & d \end{pmatrix}t\). We claim that this is the inverse. We have 
\[
(ad - bc)^{-1} \begin{pmatrix}d & -b \\c & a \end{pmatrix} \begin{pmatrix}a & b \\c & d \end{pmatrix} = (ad - bc)^{-1} \begin{pmatrix}ad - bc & 0 \\0 & ad - bc \end{pmatrix}
\]
\[
= \begin{pmatrix}1 & 0 \\0 & 1 \end{pmatrix}.
\]

Similarly, one checks that \(\begin{pmatrix}a & b \\c & d \end{pmatrix}(ad - bc)^{-1} \begin{pmatrix}d & -b \\c & a \end{pmatrix} = \begin{pmatrix}1 & 0 \\0 & 1 \end{pmatrix}\).

Suppose now that \(\begin{pmatrix}a & b \\c & d \end{pmatrix}\) is invertible. The expression \(ad - bc\) is called the determinant of the matrix \(M = \begin{pmatrix}a & b \\c & d \end{pmatrix}\) and is denoted \(\text{det}(M)\). One can verify by a laborious but straightforward calculation that for any two matrices \(M, N\) we have
\[
\text{det}(MN) = \text{det}(M) \text{det}(N).
\]
If the matrix \(M\) has an inverse, say \(MN = NM = I_2\), then
\[
\text{det}(MN) = \text{det}(M) \text{det}(N) = \text{det}(I_2) = 1,
\]
and that shows that \(\text{det}(M) \neq 0\). (One can then show that \(N\) is necessarily \((ad - bc)^{-1} \begin{pmatrix}d & -b \\c & a \end{pmatrix}\), but we don’t really need this fact now).

19. Quotient rings

Consider a surjective ring homomorphism \(f : R \rightarrow S\). Given an element \(s \in S\) let \(r\) be an element of \(R\) such that \(f(r) = s\). How unique is \(r\)? If \(a \in I := \text{Ker}(f)\) then \(f(r + a) = f(r) + f(a) = f(r) + 0 = f(r)\). Conversely, if \(f(r_1) = s\) then \(f(r_1 - r) = f(r_1) - f(r) = s - s = 0\) so \(a := r_1 - r \in I\) and \(r_1 = r + a\). Let us use the notation
\[
r + I = \{r + i : i \in I\}
\]
(this called a coset of \(I\)). Then we have proven that
\[
f^{-1}(s) = r + I.
\]
Thus, in a sense, we may identify elements of \(S\) with cosets of \(R\) and from this point of view we may that the cosets (thought of as being the elements of \(S\)) form a ring.

In this section we perform a key construction that eliminates the need in \(S\). Given a ring \(R\) and a two sided ideal \(I \triangleleft R\) we construct a new ring \(R/I\), whose elements are cosets of \(I\).
Definition 19.1. Let $R$ be a ring and $I \triangleleft R$ a two-sided ideal. A coset of $I$ is a subset of $R$ of the form

$$a + I := \{a + i : i \in I\},$$

where is $a$ is an element of $R$.

Example 19.2. Suppose that $R = \mathbb{Z}$ and $I = (n)$ for some positive integer $n$. Then $a + (n) = \{\ldots, a - n, a, a + n, a + 2n, \ldots\}$ are precisely the numbers congruent to $a$ modulo $n$.

Lemma 19.3. We have the following facts:

1. Every element of $R$ belongs to a coset of $I$.
2. Two cosets are either equal or disjoint.
3. The following are equivalent: (i) $a + I = b + I$; (ii) $a \in b + I$; (iii) $a - b \in I$.

Proof. The first claim is easy: the element $r$ belongs to the coset $r + I$, because $r = r + 0$ and $0 \in I$.

Suppose that $a + I \cap b + I \neq \emptyset$. Then, there is an element of $R$ that can be written as

$$a + i_1 = b + i_2,$$

for some $i_1, i_2 \in I$. We show that $a + I \subset b + I$; by symmetry we have the opposite inclusion and so the cosets are equal. An element of $a + I$ has the form $a + i$ for some $i \in I$. We have $a + i = b + (i_2 - i_1) + i = b + (i_2 - i_1 + i)$. Note that $i_2 - i_1 + i \in I$ and so $a + i \in b + I$.

We next prove the equivalence of (i), (ii) and (iii). Clearly (i) implies (ii) because $a = a + 0 \in a + I$. If (ii) holds then $a = b + i$ for some $i \in I$ and so $a - b = i \in I$ and (iii) holds. If (iii) holds then $a - b = i$ for some $i \in I$ and so $a = a + 0 \in a + I$, but also $a = b + i \in b + I$. That is, $a + I \cap b + I \neq \emptyset$ and so $a + I = b + I$.

Theorem 19.4. Let $R$ be a ring and $I \triangleleft R$ a two-sided ideal. Denote the collection of cosets of $I$ in $R$ by $R/I$. Define addition by

$$(a + I) + (b + I) = a + b + I,$$

and multiplication by

$$(a + I)(b + I) = ab + I.$$ 

These operations are well defined and make $R/I$ into a ring (a quotient ring) with zero element $0 + I = I$ and identity element $1 + I$.

Proof. First, our definition of the operations makes use of writing a coset as $a + I$. This way of writing is not unique and so we should check that our definitions are independent of the choice of $a$ such that the coset is $a + I$. Namely, if

$$a + I = a' + I, b + I = b' + I,$$

we need to check that

$$a + b + I = a' + b' + I, \quad ab + I = a'b' + I.$$ 

Now, $(a + b) - (a' + b') = (a - a') + (b - b')$. By the lemma above $a - a' \in I, b - b' \in I$ and so $(a + b) - (a' + b') \in I$ and so, by the same lemma, $a + b + I = a' + b' + I$. Also $ab - a'b' =
\( (a - a')b + a'(b - b') \). Now, \( a - a' \in I, b - b' \in I \) and so \((a - a')b \in I, a'(b - b') \in I \) and it follows that \( ab - a'b' \in I \). Therefore, \( ab + I = a'b' + I \).

We now verify the ring axioms. It will be convenient to write \( \bar{a} \) for \( a + I \). With this notation we have

\[
\bar{a} + \bar{b} = \bar{a + b}, \quad \bar{a} \cdot \bar{b} = \bar{ab}.
\]

The axioms follow from the definition of the operations and the fact that they hold for \( R \). To make clear when is it that we use that the axioms hold in \( R \) we use \( \equiv \) at that point.

1. \( \bar{a} + \bar{b} = \bar{a + b} \equiv \bar{b + a} = \bar{b + \bar{a}}. \)
2. \( \bar{a} + (\bar{b} + \bar{c}) \equiv \bar{a + (b + c)} = \bar{a + b} + \bar{c} = (\bar{a + b}) + \bar{c}. \)
3. \( \bar{0} + \bar{a} = \bar{0 + a} \equiv \bar{a}. \) (We remark that \( \bar{0} = \bar{I} \).)
4. \( \bar{a} + \bar{-a} = \bar{a + (-a)} \equiv \bar{0}. \)
5. \( \bar{a}\bar{b} \bar{c} = \bar{a\cdot bc} = \bar{a(bc)} \equiv (\bar{a}b)c = \bar{ab}\bar{c} = (\bar{a} \bar{b}) \bar{c}. \)
6. \( \bar{a} \bar{1} = \bar{a} \equiv \bar{1} \bar{a} = \bar{1a} \equiv \bar{a}. \)
7. \( \bar{a} + \bar{b} \bar{c} = \bar{a + b} \bar{c} = \bar{(a + b)c} \equiv \bar{ac + bc} = \bar{ac} + \bar{bc} = \bar{a} \bar{c} + \bar{b} \bar{c}. \) Also, \( \bar{c}(\bar{a} + \bar{b}) = \bar{c} \bar{a} + \bar{c} \bar{b} = \bar{c(a + b)} \equiv \bar{ca + cb} = \bar{ca} + \bar{cb} = \bar{c} \bar{a} + \bar{c} \bar{b}. \)

\[ \square \]

**Proposition 19.5.** The natural map,

\[ \pi : R \to R/I, \quad a \mapsto \pi(a) := \bar{a} \]

is a surjective ring homomorphism with kernel \( I \). Thus, every ideal \( I \lhd R \) is the kernel of some ring homomorphism from \( R \) to some other ring.

**Proof.** Note that \( 1 \mapsto \bar{1} \), which is the identity element of \( R/I \). We have \( \pi(a + b) = \bar{a + b} = \bar{a + \bar{b}} = \pi(a) + \pi(b) \). Also, \( \pi(ab) = \bar{ab} = \bar{a} \bar{b} = \pi(a) \pi(b) \). We have shown that \( \pi \) is a ring homomorphism and it is clearly surjective.

The kernel of \( \pi \) are the elements \( a \in R \) such that \( a + I = 0 + I \). By the lemma above this is the set of elements \( a \) such that \( a - 0 \in I \), namely, the kernel is precisely \( I \). \[ \square \]

**Example 19.6.** Consider the ring \( \mathbb{Z} \). If we take the ideal \( \{0\} \) then \( \mathbb{Z}/\{0\} \) can be identified with \( \mathbb{Z} \); the map \( \mathbb{Z} \to \mathbb{Z}/\{0\} \) is a bijective ring homomorphism. Let \( n > 0 \) then. The ring \( \mathbb{Z}/(n) \) has as elements the cosets \( a + (n) \). Two cosets \( a + (n), b + (n) \) are equal if and only if \( a - b \in (n) \), that is, precisely when \( n|(a-b) \). We see that the elements of \( \mathbb{Z}/(n) \) are just the congruence classes modulo \( n \) and the operations on \( \mathbb{Z}/(n) \) are then just the operations we defined on congruence classes.
19.1. **The quotient ring** $\mathbb{F}[x]/(f(x))$. Let $\mathbb{F}$ be a field, $f(x) \in \mathbb{F}[x]$ a non-constant polynomial. Consider the quotient ring $\mathbb{F}[x]/(f(x))$. Suppose that $f(x) = x^n + a_{n-1}x^{n-1} + \cdots + a_0$ is a monic polynomial of degree $n$. The following lemma is an analogue of Lemma 12.1.

**Lemma 19.7.** Every element of $\mathbb{F}[x]/(f(x))$ is of the form $\overline{g(x)} = g(x) + (f(x))$ for a unique polynomial $g(x)$ which is either zero or of degree less than $n$.

**Proof.** Let $h(x)$ be a polynomial. To say that that $h(x) + (f(x)) = g(x) + (f(x))$ is to say that $h(x) = q(x)f(x) + g(x)$. The requirement that $\deg(g) < \deg(f)$ is exactly to say that the expression

$$h(x) = q(x)f(x) + g(x)$$

is dividing $h$ by $f$ with residue. We know this is always possible and in a unique fashion. \(\square\)

**Theorem 19.8.** Let $\mathbb{F}$ be a field, $f(x) \in \mathbb{F}[x]$ a non-constant irreducible polynomial of degree $n$. Then $\mathbb{F}[x]/(f(x))$ is a field. If $\mathbb{F}$ is a finite field of cardinality $q$ then $\mathbb{F}[x]/(f(x))$ is a field of $q^n$ elements.

**Proof.** We already know that $\mathbb{F}[x]/(f(x))$ is a commutative ring. We note that $\overline{1} \neq \overline{0}$ because $1 = 1 - 0 \notin (f)$ (if it did, $f$ would be a constant polynomial). Thus, we only need to show that a non-zero element has an inverse. Let $\overline{g(x)}$ be a non-zero element. That means that $f(x) \nmid g(x)$ and so that $\gcd(f, g) = 1$ (here is where we use that $f$ is irreducible). Therefore, there are polynomials $u(x), v(x)$ such that

$$u(x)f(x) + v(x)g(x) = 1.$$ 

Passing to the quotient ring, that means that $\overline{\bar{v}g} = \overline{1}$, which is the identity of the quotient ring.

Finally, by the Lemma, every element of $\mathbb{F}[x]/(f(x))$ has a unique representative of the form $a_{n-1}x^{n-1} + \cdots + a_1x + a_0$. If $\mathbb{F}$ has $q$ elements, we get $q^n$ such polynomials. \(\square\)

**Example 19.9. A field with 4 elements.** Take the field $\mathbb{F}$ to be $\mathbb{Z}/2\mathbb{Z}$ and consider the polynomial $x^2 + x + 1$ over that field. Because it is of degree 2 and has no root in $\mathbb{Z}/2\mathbb{Z}$ it must be irreducible. Therefore, $\mathbb{Z}/2\mathbb{Z}[x]/(x^2 + x + 1)$ is a field $\mathbb{K}$ with 4 elements. Let us list its elements:

$$\mathbb{K} = \{0, \overline{1}, \overline{x}, \overline{x+1}\}.$$ 

We can give tables of addition and multiplication:

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Example 19.10. A field with 9 elements. Consider the polynomial \( x^2 + 1 \) over \( \mathbb{Z}/3\mathbb{Z} \). It is quadratic and has no root in \( \mathbb{Z}/3\mathbb{Z} \), hence is irreducible over \( \mathbb{Z}/3\mathbb{Z} \). We conclude that \( L = \mathbb{Z}/3\mathbb{Z}[x]/(x^2 + 1) \) is a field with 9 elements. Note that in \( \mathbb{Z}/3\mathbb{Z} \) the element \(-1 = 2\) is not a square. However, in \( L \) we have \( x^2 = x^2 - (x^2 + 1) = -1 \) and so \(-1\) is a square now – its root is \( x \) (viewed as an element of \( L \)). In fact, any quadratic polynomial over \( \mathbb{Z}/3\mathbb{Z} \) has a root in \( L \), because the discriminant \( b^2 - 4ac \) is either 0, 1, 2 and all those are squares in \( L \).

For example, consider the polynomial \( t^2 + t + 2 \). It has discriminant \(-7 \equiv -1 \pmod{3}\), which is not a square in \( \mathbb{Z}/3\mathbb{Z} \). In the field \( L \) we have \( x^2 = -1 \). The solutions of the polynomial are then \( (-1 \pm \sqrt{-1})/2 = 2(-1 \pm x) = 1 \pm 2x \).

On the other hand, one can prove that the polynomial \( t^3 + t^2 + 2 \) is irreducible in \( \mathbb{Z}/3\mathbb{Z} \) and stays irreducible in \( L \). In MATH 370 we learn a systematic theory for deciding which polynomials stay irreducible and which not.

Example 19.11. Fields with 8 and 16 elements. A polynomial of degree 3 is irreducible if and only if it doesn’t have a root. We can verify that \( x^3 + x + 1 \) doesn’t have a root in \( \mathbb{Z}/2\mathbb{Z} \) and conclude that \( \mathbb{Z}/2\mathbb{Z}[x]/(x^3 + x + 1) \) is a field with 8 elements. Consider the field \( K \) with 4 elements constructed above. We note that the polynomial \( t^2 + t + \bar{x} \) is irreducible over \( K \) (simply by substituting for \( t \) any of the four elements of \( K \) and checking). Thus, we get a field \( L \) with 16 elements

\[ L = K[t]/(t^2 + t + \bar{x}). \]

19.2. Every polynomial has a root in a bigger field.

**Theorem 19.12.** Let \( F \) be a field and \( f(x) \in F[x] \) a non-constant polynomial. There is a field \( L \) containing \( F \) and an element \( \ell \in L \) such that \( f(\ell) = 0 \).

**Proof.** If \( g|f \) and \( g(\ell) = 0 \) then also \( f(\ell) = 0 \), so we may assume that \( f \) is irreducible. Let \( L = F[x]/(f(x)) \). This is a field. We have a natural map \( F \to L, a \mapsto \bar{a} \). This map is an injective ring homomorphism and we identify \( F \) with its image in \( L \) so as to say that \( L \supset F \).

Now, suppose that \( f(x) = a_nx^n + \ldots + a_1x + a_0 \). To say that \( f \) has a root in \( L \) is to say that for some element \( \ell \in L \) we have

\[ a_n\ell^n + \ldots + a_1\ell + a_0 = 0. \]

We check that this hold for the element \( \ell = \bar{x} \). Indeed,

\[ a_n\bar{x}^n + \ldots + a_1\bar{x} + a_0 = a_nx^n + \ldots + a_1x + a_0 = f(x) = 0_L. \]

Example 19.13. According to this result, \(-1\) has a square root in the field \( \mathbb{R}[x]/(x^2 + 1) \). One can show that \( \mathbb{R}[x]/(x^2 + 1) \cong \mathbb{C} \).
19.3. Roots of polynomials over \( \mathbb{Z}/p\mathbb{Z} \). We can now continue our discussion, begun in §15.7 of the efficient determination of whether a small degree polynomial \( f(x) \) over \( \mathbb{Z}/p\mathbb{Z} \) has a root in \( \mathbb{Z}/p\mathbb{Z} \). Recall that the only remaining point was whether the Euclidean algorithm step,

\[ x^p - x = q(x)f(x) + r(x), \]

can be done rapidly. Now we can answer that affirmatively. Note that \( r(x) + x \) is exactly the representative of \( x^p \) in the ring \( \mathbb{F}[x]/(f(x)) \). This representative can be calculated quickly by the method we already used for calculating powers. We need to calculate

\[ x, x^2, x^4, x^8, \ldots \]

and the express \( p \) in base 2, \( p = \sum a_i2^i, a_i \in \{0,1\} \), \( x^p = \prod_{\{i,a_i \neq 0\}} x^{2^i} \) and so. We see that the slowing factor now is how quickly we can carry out multiplication in the ring \( \mathbb{F}[x]/(f(x)) \). It is not hard to see that this depends on the degree of \( f \) and not on \( p \).

20. The First Isomorphism Theorem

20.1. Isomorphism of rings.

**Definition 20.1.** Let \( R, S \) be rings. A ring homomorphism \( f : R \to S \) is called an isomorphism if \( f \) is bijective.

**Lemma 20.2.** If \( f : R \to S \) is a ring isomorphism then the inverse function \( g = f^{-1}, g : S \to R \) is also a ring homomorphism, hence an isomorphism. (The inverse function is defined by \( g(s) = r \), where \( r \) is the unique element such that \( f(r) = s \).)

**Proof.** First, because \( f(1_R) = 1_S \) we have \( g(1_S) = 1_R \). Next, let \( s_1, s_2 \in S \). We need to prove \( g(s_1 + s_2) = g(s_1) + g(s_2) \) and \( g(s_1s_2) = (g(s_1))g(s_2) \). It is enough to prove that

\[
 f(g(s_1 + s_2)) = f(g(s_1) + g(s_2)),
 \quad f(g(s_1s_2)) = f(g(s_1))g(s_2)),
\]

because \( f \) is injective. But \( f(g(s_1) + g(s_2)) = f(g(s_1)) + f(g(s_2)) = s_1 + s_2 = f(g(s_1 + s_2)) \) and \( f(g(s_1)g(s_2)) = f(g(s_1))f(g(s_2)) = s_1s_2 = f(g(s_1))f(g(s_2)) \).

**Definition 20.3.** Let \( R, S \) be rings. We say that \( R \) and \( S \) are isomorphic if there is a ring isomorphism \( R \to S \).

**Lemma 20.4.** Being isomorphic is an equivalence relation on rings.

**Proof.** First, the identity function is always a ring homomorphism from \( R \) to \( R \), so this relation is reflexive. Secondly, if \( f : R \to S \) is an isomorphism then \( g : S \to R \) is an isomorphism, where \( g \) is the inverse function to \( f \). Thus, the relation is symmetric. Now suppose \( f : R \to S \) and \( g : S \to T \) are ring isomorphisms between the rings \( R, S, T \). To show the relation is transitive we need to prove that \( g \circ f : R \to T \) is an isomorphism. Indeed:

1. \( (g \circ f)(1_R) = g(f(1_R)) = g(1_S) = 1_T; \)
(2) \((g \circ f)(r_1 + r_2) = g(f(r_1 + r_2)) = g(f(r_1) + f(r_2)) = (g \circ f)(r_1) + (g \circ f)(r_2)\);

(3) \((g \circ f)(r_1 r_2) = g(f(r_1 f(r_2))) = g(f(r_1))g(f(r_2)) = (g \circ f)(r_1)(g \circ f)(r_2)\). 

We shall denote the fact that \(R\) is isomorphic to \(S\) by \(R \cong S\).

20.2. The First Isomorphism Theorem.

**Theorem 20.5.** Let \(f : R \rightarrow S\) be a surjective homomorphism of rings. Let \(I = \ker(f)\) then there is an isomorphism \(F : R/I \rightarrow S\), such that the following diagram commutes

\[
\begin{array}{ccc}
R & \xrightarrow{f} & S \\
\downarrow{\pi} & & \downarrow{F} \\
R/I & & \\
\end{array}
\]

where \(\pi : R \rightarrow R/I\) is the canonical map \(g \mapsto \bar{g}\).

**Proof.** We define a function, \(F : R/I \rightarrow S\), by \(F(\bar{g}) = f(g)\). We first prove that this map is well defined. Suppose that \(\bar{g} = \bar{g_1}\). We need to show that \(f(g) = f(g_1)\). This holds because \(\bar{g} = \bar{g_1}\) means \(g - g_1 \in I = \ker(f)\). Now:

- \(F(1_{R/I}) = F(\bar{1}_R) = f(1_R) = 1_S;\)
- \(F(\bar{g} + \bar{h}) = F(\bar{g} + \bar{h}) = f(g + h) = f(g) + f(h) = F(\bar{g}) + F(\bar{h});\)
- \(F(\bar{g} \cdot \bar{h}) = F(gh) = f(g) \cdot f(h) = F(\bar{g}) \cdot F(\bar{h}).\)

We also have

\((F \circ \pi)(g) = F(\bar{g}) = f(g),\)

so \(F \circ \pi = f\). Because of this we have that \(F\) is surjective. We next show \(F\) is injective. Suppose that \(F(\bar{g}) = 0_S\) then \(f(g) = 0_S\) and so \(g \in I\). Thus, \(\bar{g} = 0_R/I\). □

**Example 20.6.** We consider again the homomorphism \(\mathbb{Z} \rightarrow \mathbb{Z}/n\mathbb{Z}\). It is a surjective ring homomorphism with kernel \((n)\) and we conclude that \(\mathbb{Z}/(n) \cong \mathbb{Z}/n\mathbb{Z}\), a fact we have noticed somewhat informally before.

**Example 20.7.** We have \(\mathbb{R}[x]/(x^2 + 1) \cong \mathbb{C}\). To show that, define a ring homomorphism

\[\mathbb{R}[x] \rightarrow \mathbb{C},\]

by \(\sum_{j=0}^{n} a_j x^j \mapsto \sum_{j=0}^{n} a_j i^j\). This is a well defined function taking 1 to 1. It is easy to verify it is a homomorphism. In fact, recall that \(\mathbb{C}[x] \rightarrow \mathbb{C}, f \mapsto f(i)\), is a homomorphism. Our
map is the restriction of the evaluation-at-i homomorphism to the subring $\mathbb{R}[x]$ and so is also a homomorphism. It is also clear that this homomorphism is surjective. The kernel $I$ definitely contains $x^2 + 1$, and so all its multiples. That is, $I$ contains the ideal $(x^2 + 1)$. Now, $I = (f)$ for some polynomial $f$. Because $x^2 + 1 \in (f)$, $f|(x^2 + 1)$. Since $x^2 + 1$ is irreducible over $\mathbb{R}$, either $f \sim 1$ or $f \sim x^2 + 1$. If $f \sim 1$ we have $(f) = \mathbb{R}[x]$ and so any polynomial is in the kernel, which is clearly not the case (for example, 1 is not in the kernel). Thus $f \sim x^2 + 1$ and $I = (x^2 + 1)$; by the first isomorphism theorem we have

$$\mathbb{R}[x]/(x^2 + 1) \cong \mathbb{C}.$$  

20.3. The Chinese Remainder Theorem.

**Theorem 20.8.** Let $m, n$ be positive integers such that $(m, n) = 1$. Then

$$\mathbb{Z}/mn\mathbb{Z} \cong \mathbb{Z}/m\mathbb{Z} \times \mathbb{Z}/n\mathbb{Z}.$$  

**Proof.** We define a function

$$f : \mathbb{Z} \to \mathbb{Z}/m\mathbb{Z} \times \mathbb{Z}/n\mathbb{Z}, \quad f(a) = (a \pmod{m}, a \pmod{n}).$$  

This function is a ring homomorphism:

- $f(1) = (1 \pmod{m}, 1 \pmod{n}) = (1_{\mathbb{Z}/m\mathbb{Z}}, 1_{\mathbb{Z}/n\mathbb{Z}});$  
- $f(a + b) = (a + b \pmod{m}, a + b \pmod{n}) = (a \pmod{m}, a \pmod{n}) + (b \pmod{m}, b \pmod{n}) = f(a) + f(b);$  
- $f(ab) = (ab \pmod{m}, ab \pmod{n}) = (a \pmod{m}, a \pmod{n}) \cdot (b \pmod{m}, b \pmod{n}) = f(a)f(b).$

The kernel of the map is the set $\{a : m|a, n|a\} = \{a : mn|a\}$ (using that $(m, n) = 1$), that is, the kernel is the principal ideal $(mn)$. That means that the integers $0, 1, \ldots, mn - 1$ all have different images in the target. Since the target has $mn$ elements, we conclude that $f$ is surjective. By the first isomorphism theorem

$$\mathbb{Z}/(mn) \cong \mathbb{Z}/m\mathbb{Z} \times \mathbb{Z}/n\mathbb{Z}.$$  

□  

This theorem is very useful. It says that to solve an equation modulo $mn$, $(m, n) = 1$, is the same as solving it modulo $m$ and modulo $n$. That is, for given integers $a_0, \ldots, a_n$ and an integer $A$ we have $a_n A^n + \cdots + a_1 A + a_0 \equiv 0 \pmod{mn}$ if and only if we have $a_n A^n + \cdots + a_1 A + a_0 \equiv 0 \pmod{m}$ and $a_n A^n + \cdots + a_1 A + a_0 \equiv 0 \pmod{n}$. Here is an example:

**Example 20.9.** Solve the equation $5x + 2 = 0$ modulo 6.

We consider the equation modulo 2 and get $x = 0 \pmod{2}$; we consider it modulo 3 and get $2x + 2 = 0 \pmod{3}$ and get that $x = 2 \pmod{3}$. There is an $x \in \mathbb{Z}$ such that $x \pmod{2} = 0, x \pmod{3} = 2$. Therefore, $x = 2 \pmod{6}$. So we have $5x + 2 = 0 \pmod{6}$.
\((\text{mod } 3) = 2\) and in fact \(x\) is unique modulo 6 (this is the CRT). We can guess that \(x = 2\) will do in this case, but it raises the general problem of finding the inverse isomorphism to
\[ \mathbb{Z}/mn\mathbb{Z} \to \mathbb{Z}/m\mathbb{Z} \times \mathbb{Z}/n\mathbb{Z}. \]

20.3.1. Inverting \(\mathbb{Z}/mn\mathbb{Z} \to \mathbb{Z}/m\mathbb{Z} \times \mathbb{Z}/n\mathbb{Z}\). Suppose we know how to find elements \(e_1, e_2\) such that \(e_1 = 1 \pmod{m}, e_1 = 0 \pmod{n}\) and \(e_2 = 0 \pmod{m}, e_2 = 1 \pmod{n}\), then we would have solved our problem. Indeed, given now two congruence classes \(a \pmod{m}, b \pmod{n}\) take the element \(ae_1 + be_2\).

Example 20.10. Solve the equation \(56x + 23 = 0 \pmod{323}\).
We have \(323 = 17 \cdot 19\).
- **Solution modulo 17.**
  We have the equation \(5x + 6 = 0 \pmod{17}\). Or \(x = -6 \cdot 5^{-1} = 11 \cdot 5^{-1}\). To find \(5^{-1}\) we look for \(u, v\) such that \(1 = u5 + v17\).
  
  \(17 = 3 \cdot 5 + 2, \, 5 = 2 \cdot 2 + 1\) so \(1 = 5 - 2 \cdot 2 = 5 - 2 \cdot (17 - 3 \cdot 5) = 7 \cdot 5 - 2 \cdot 17\) and so \(7 \cdot 5 = 1 \pmod{17}\). We conclude that \(x = 11 \cdot 7 = 77 = 9 \pmod{17}\).
- **Solution modulo 19.**
  We have the equation \(-x + 4 = 0\) so \(x = 4 \pmod{19}\) is a solution.
- **Finding \(e_1, e_2\).**
  We have \(19 = 17 + 2, \, 17 = 8 \cdot 2 + 1\) so \(1 = 17 - 8 \cdot 2 = 9 \cdot 17 - 8 \cdot 19\). It follows that \(e_1 = 1 - 9 \cdot 17 = -152, e_2 = 1 + 8 \cdot 19 = 153\).
- We conclude that the solution to the equation \(56x + 23 = 0 \pmod{323}\) is \(9 \cdot e_1 + 4 \cdot e_2 = -1368 + 612 = -756\) and modulo 323 this is 213.

Example 20.11. Solve the equation \(x^2 = 118 \pmod{323}\). As before we reduce to solving \(x^2 = 118 = 16 \pmod{17}\) and \(x^2 = 118 = 4 \pmod{19}\). There are two solutions in each case, given by \(x = \pm 4 \pmod{17}\) and \(x = \pm 2 \pmod{19}\). We conclude that over all we have 4 solutions given by \(\pm 4(-152) \pm 2(153) \pmod{323}\).

One can reduce those numbers to standard representatives and find that \(21, 55, 268, 302 \pmod{323}\) are the four solutions.

One can be more precise about the connection between solutions \(\text{mod } mn\) and solutions \(\text{mod } m\) and \(\text{mod } n\). First, let us generalize the Chinese Remainder Theorem:

**Theorem 20.12.** Let \(m_1, \ldots, m_k\) be relatively prime non-zero integers (that is \((m_i, m_j) = 1\) for \(i \neq j\)). Then there is an isomorphism
\[ \mathbb{Z}/m_1m_2\ldots m_k\mathbb{Z} \cong \mathbb{Z}/m_1\mathbb{Z} \times \mathbb{Z}/m_2\mathbb{Z} \times \cdots \times \mathbb{Z}/m_k\mathbb{Z}, \]
given by
\[ a \pmod{m_1m_2\ldots m_k} \mapsto (a \pmod{m_1}, a \pmod{m_2}, \ldots, a \pmod{m_k}). \]
The theorem is not hard to prove by induction on $k$. The main case, $k = 2$, is the one we proved above. Now, let $g(x) = a_nx^n + \ldots a_1x + a_0$ be a polynomial with integer coefficients. Let $S$ be the solutions of $g$ in $\mathbb{Z}/m_1m_2\ldots m_k\mathbb{Z}$ and $S_i$ the solutions of $g$ in $\mathbb{Z}/m_i\mathbb{Z}$. Then we have a bijection

$$S \leftrightarrow S_1 \times S_2 \times \cdots \times S_k,$$

given by

$$a \pmod{m_1m_2\ldots m_k} \mapsto (a \pmod{m_1}, a \pmod{m_2}, \ldots, a \pmod{m_k}).$$

Indeed, the expression $g(a) \pmod{m_1m_2\ldots m_k}$ is mapped to

$$(g(a) \pmod{m_1}, g(a) \pmod{m_2}, \ldots, g(a) \pmod{m_k})$$

and $g(a) \equiv 0 \pmod{m_1m_2\ldots m_k}$ if and only if for every $i$ we have $g(a) \equiv 0 \pmod{m_i}$. This shows that we have a map

$$S \rightarrow S_1 \times S_2 \times \cdots \times S_k.$$

But, conversely, given solutions $r_i$ to $g(x) \pmod{m_i}$, there is a unique $r \pmod{m_1m_2\ldots m_k}$ such that $r \equiv r_i \pmod{m_i}$ and $g(r) \equiv 0 \pmod{m_1m_2\ldots m_k}$ because it is true modulo every $m_i$.

In particular, we may draw the following conclusion.

**Corollary 20.13.** Let $m_1, \ldots, m_k$ be relatively prime integers. Let $s$ be the number of solutions to the equation $a_nx^n + \ldots a_1x + a_0 = 0 \pmod{m_1m_2\ldots m_k}$ and let $s_i$ be the number of solutions modulo $m_i$. Then

$$s = s_1s_2\cdots s_k.$$

**Example 20.14.** The equation $x^2 = 1$ has 8 solutions modulo $2 \cdot 3 \cdot 5 \cdot 7 = 490$, because it has one solution mod 2, and 2 solutions mod 3,5 or 7.

**Example 20.15.** The equation $34x = 85 \pmod{17 \cdot 19}$ has 17 solutions, because it has 17 solutions modulo 17 (it is then the equation $0 \cdot x = 0 \pmod{17}$) and has a unique solution modulo 19 (it is then the equation $4x = 10 \pmod{19}$ and $x = 12$ is the unique solution).

On the other hand, the equation $34x = 5 \pmod{17 \cdot 19}$ has no solutions, because it has no solutions modulo 17.

There remains the question how to calculate a solution mod $m_1m_2\ldots m_k$ from solutions mod $m_1$, mod $m_2$, \ldots, mod $m_k$. That is, how to find explicitly the inverse to the map

$$\mathbb{Z}/m_1m_2\ldots m_k\mathbb{Z} \rightarrow \mathbb{Z}/m_1\mathbb{Z} \times \mathbb{Z}/m_2\mathbb{Z} \times \cdots \times \mathbb{Z}/m_k\mathbb{Z}.$$ 

We explain how to do that for 3 numbers $m_1, m_2, m_3$, though the method is general.

We first find integers $\epsilon_1, \epsilon_2$ such that

$$\epsilon_1 \equiv 1 \pmod{m_1}, \quad \epsilon_1 \equiv 0 \pmod{m_2m_3},$$

and

$$\epsilon_2 \equiv 0 \pmod{m_1}, \quad \epsilon_2 \equiv 1 \pmod{m_2m_3}.$$
This we know how to do because we are only dealing with two relatively prime numbers, that is, \( m_1 \) and \( m_2 m_3 \). Then, find \( \lambda_2, \lambda_3 \) such that
\[
\lambda_2 \equiv 1 \pmod{m_2}, \quad \lambda_2 \equiv 0 \pmod{m_3},
\]
and
\[
\lambda_3 \equiv 0 \pmod{m_2}, \quad \lambda_3 \equiv 1 \pmod{m_3}.
\]
Then, the numbers

\[
\mu_1 = \epsilon_1, \quad \mu_2 = \epsilon_2 \lambda_2, \quad \mu_3 = \epsilon_2 \lambda_3,
\]
are congruent to \((1,0,0),(0,1,0)\) and \((0,0,1)\) respectively in the ring \( \mathbb{Z}/m_1 \mathbb{Z} \times \mathbb{Z}/m_2 \mathbb{Z} \times \mathbb{Z}/m_3 \mathbb{Z} \).

To find an integer \( \mod{m_1 m_2 m_3} \) mapping to \((a,b,c)\) in \( \mathbb{Z}/m_1 \mathbb{Z} \times \mathbb{Z}/m_2 \mathbb{Z} \times \mathbb{Z}/m_3 \mathbb{Z} \) take \( a \mu_1 + b \mu_2 + c \mu_3 \):

\[
a \mu_1 + b \mu_2 + c \mu_3 \mapsto (a,b,c) \in \mathbb{Z}/m_1 \mathbb{Z} \times \mathbb{Z}/m_2 \mathbb{Z} \times \mathbb{Z}/m_3 \mathbb{Z}.
\]

Let us illustrate all this with a numerical example.

**Example 20.16.** Find the solutions to the equation \( x^2 + x + 2 = 0 \pmod{7 \cdot 11 \cdot 23} \).

The solutions modulo 7 are \( S_1 = \{3\} \); the solutions modulo 11 are \( S_2 = \{4, 6\} \); the solutions modulo 23 are \( S_3 = \{9, 13\} \). (Those are found by brute computation.) To find the corresponding solutions modulo \( 7 \cdot 11 \cdot 23 \) we first find the \( \mu_i \) above.

First \( 11 \cdot 23 = 36 \cdot 7 + 1 \) and so \( \epsilon_1 = 253, \epsilon_2 = -252 \) satisfy \( \epsilon_1 \equiv 1 \pmod{7}, \epsilon_1 \equiv 0 \pmod{253}, \)
\( \epsilon_2 \equiv 0 \pmod{7}, \epsilon_2 \equiv 1 \pmod{253} \). To find \( \lambda_1, \lambda_2 \) we note that \( 1 = 23 - 2 \cdot 11 \) and so \( \lambda_1 = 23, \lambda_2 = -22 \) satisfy \( \lambda_1 \equiv 1 \pmod{11}, \lambda_1 \equiv 0 \pmod{23} \) and \( \lambda_2 \equiv 0 \pmod{11}, \lambda_2 \equiv 1 \pmod{23} \). The numbers \( \mu_1, \mu_2, \mu_3 \) are then \( 253, -252 \cdot 23, -252 \cdot (-22) \). They should be understood as numbers modulo \( 7 \cdot 11 \cdot 23 = 1771 \) and so we can replace \( -252 \cdot 23 \) by 1288 (which is congruent to it modulo 1771) and \( -252 \cdot (-22) \) by 231. Then to get a number congruent to \((a,b,c)\) in \( \mathbb{Z}/7 \mathbb{Z} \times \mathbb{Z}/11 \mathbb{Z} \times \mathbb{Z}/23 \mathbb{Z} \), we take \( a \cdot 253 + b \cdot 1288 + c \cdot 231 \). In particular, the solution \((3,4,9)\) produces \( 3 \cdot 253 + 4 \cdot 1288 + 9 \cdot 231 = 7990 \) which we can replace by the number 906 which is congruent to it modulo 1771. By our theory 906 is one of the 4 solutions to the equation \( x^2 + x + 2 = 0 \pmod{1771} \) and this can be verified.

Do you think you could have found this solution in an easier way?! (I can’t see how).

21. **Prime and maximal ideals**

(planned)
Part 6. Groups

22. First definitions and examples

22.1. Definitions and some formal consequences.

Definition 22.1. A group $G$ is a non-empty set with an operation

$$G \times G \to G, \quad (a, b) \mapsto ab,$$

such that the following axioms hold:

1. $(ab)c = a(bc)$. (Associativity)
2. There exists an element $e \in G$ such that $eg = ge$ for all $g \in G$. (Identity)
3. For every $g \in G$ there exists an element $d \in G$ such that $dg = gd = e$. (Inverse)

Here are some formal consequences of the definition:

1. $e$ is unique. Say $\tilde{e}$ has the same property then $\tilde{e} = e\tilde{e}$, using the property of $e$, but also $e\tilde{e} = e$. Thus, $e = \tilde{e}$.
2. $d$ appearing in (3) is unique (therefore we shall call it “the inverse of $g$” and denote it by $g^{-1}$). Say $\tilde{d}$ also satisfies $\tilde{d}g = gd = e$. Then

$$\tilde{d} = d\tilde{e} = d(gd) = (dg)d = ed = d.$$
3. Cancellation: $ab = cb \Rightarrow a = c$, and $ba = bc \Rightarrow a = c$.
   If $ab = cb$ then $(ab)b^{-1} = (cb)b^{-1}$ and so $a = a(bb^{-1}) = c(bb^{-1}) = c$.
4. $(ab)^{-1} = b^{-1}a^{-1}$.
   To show that we need to show that $b^{-1}a^{-1}$ “functions as the inverse of $ab$”. We have $(ab)(b^{-1}a^{-1}) = a(bb^{-1})a^{-1} = aea^{-1} = aa^{-1} = e$. Similarly, $(b^{-1}a^{-1})(ab) = b^{-1}(a^{-1}a)b = b^{-1}eb = b^{-1}b = e$.
5. $(a^{-1})^{-1} = a$.
   This is because $aa^{-1} = a^{-1}a = e$ also shows that $a$ is the inverse of $a^{-1}$.
6. Define $a^0 = e$, $a^n = a^{n-1}a$ for $n > 0$ and $a^n = (a^{-1})^{-n}$ for $n < 0$. Then we have

$$a^ma^n = a^{m+n}, \quad (a^m)^n = a^{mn}.$$

22.2. Examples.

Example 22.2. The trivial group $G$ is a group with one element $e$ and multiplication law $ee = e$.

Example 22.3. If $R$ is a ring, then $R$ with addition only is a group. It is a commutative group. The operation in this case is of course written $g + h$. In general a group is called commutative or abelian if for all $g, h \in G$ we have $gh = hg$. It is customary in such cases to write the operation in the group as $g + h$ and not as $gh$, but this is not a must.

These examples thus include $\mathbb{Z}, \mathbb{Q}, \mathbb{R}, \mathbb{C}, F, F[\epsilon], M_2(F), \mathbb{Z}/n\mathbb{Z}$ with the addition operation.
Example 22.4. Let $R$ be a ring. Recall that the units $R^\times$ of $R$ are defined as 
\[ \{ u \in R : \exists v \in R, uv = vu = 1 \} . \]
This is a group. If $u_1, u_2 \in R$ with inverses $v_1, v_2$, respectively, then, as above, one checks that $v_2v_1$ is an inverse for $u_1u_2$ and so $R^\times$ is closed under the product operation. The associative law holds because it holds in $R$; $1_R$ serves as the identity. If $R$ is not commutative there is no reason for $R^\times$ to be commutative, though in certain cases it may be.

Thus we get the examples of $\mathbb{Z}^\times = \{ \pm 1 \}, \mathbb{Q}^\times = \mathbb{Q} - \{ 0 \}, \mathbb{R}^\times = \mathbb{R} - 0, \mathbb{C}^\times = \mathbb{C} - \{ 0 \}$, and more generally, $\mathbb{F}^\times = \mathbb{F} - \{ 0 \}$. We also have, $\text{GL}_2(F) = \{ M \in M_2(F) : \det(M) \neq 0 \}, \mathbb{F}[\epsilon]^\times = \{ a + b\epsilon : a \neq 0 \}$.

Proposition 22.5. Let $n > 1$ be an integer. The group $\mathbb{Z}/n\mathbb{Z}^\times$ is precisely 
\[ \{ 1 \leq a \leq n : (a, n) = 1 \} . \]

Proof. If $\bar{a}$ is invertible then $ab = 1$ (mod $n$) for some integer $b$; say $ab = 1 + kn$ for some $k \in \mathbb{Z}$. If $d|a, d|n$ then $d|1$. Therefore $(a, n) = 1$.

Conversely, suppose that $(a, n) = 1$ then for some $u, v$ we have $1 = ua + vn$ and so $ua = 1$ (mod $n$).

One defines Euler’s $\varphi$ function on positive integers by
\[ \varphi(n) = \begin{cases} 
1 & n = 1 \\
|\mathbb{Z}/n\mathbb{Z}^\times| & n > 1.
\end{cases} \]

One can prove that this is a multiplicative function, namely, if $(n, m) = 1$ then $\varphi(nm) = \varphi(n)\varphi(m)$.

I invite you to try and prove that based on the Chinese Remainder Theorem.

Here are some specific examples:

<table>
<thead>
<tr>
<th>$n$</th>
<th>$\mathbb{Z}/n\mathbb{Z}^\times$</th>
<th>$\varphi(n)$</th>
</tr>
</thead>
<tbody>
<tr>
<td>2</td>
<td>${1}$</td>
<td>1</td>
</tr>
<tr>
<td>3</td>
<td>${1, 2}$</td>
<td>2</td>
</tr>
<tr>
<td>4</td>
<td>${1, 3}$</td>
<td>2</td>
</tr>
<tr>
<td>5</td>
<td>${1, 2, 3, 4}$</td>
<td>4</td>
</tr>
<tr>
<td>6</td>
<td>${1, 5}$</td>
<td>2</td>
</tr>
<tr>
<td>7</td>
<td>${1, 2, 3, 4, 5, 6}$</td>
<td>6</td>
</tr>
<tr>
<td>8</td>
<td>${1, 3, 5, 7}$</td>
<td>4</td>
</tr>
<tr>
<td>9</td>
<td>${1, 2, 4, 5, 7, 8}$</td>
<td>6</td>
</tr>
</tbody>
</table>

Example 22.6. If $G, H$ are groups then $G \times H$ is a group with the operation
\[ (g_1, h_1)(g_2, h_2) = (g_1g_2, h_1h_2) . \]
The identity is \((e_G, e_H)\) and \((g, h)^{-1} = (g^{-1}, h^{-1})\).

22.3. **Subgroups.** Let \(G\) be a group. A subset \(H \subseteq G\) is called a *subgroup* if the following holds:

1. \(e_G \in H\);
2. \(a, b \in H \Rightarrow ab \in H\);
3. \(a \in H \Rightarrow a^{-1} \in H\).

Clearly then \(H\) is a group in its own right.

**Example 22.7.** The subset \(S^1\) of \(\mathbb{C}^\times\) consisting of all complex numbers of absolute value 1 is a subgroup. Indeed \(1 \in S^1\). If \(s_1, s_2 \in S^1\) then \(|s_1s_2| = |s_1|\cdot|s_2| = 1\) so \(s_1s_2 \in S^1\). If \(z\) is any non-zero complex number then \(1 = |1| = |zz^{-1}| = |z|\cdot|z^{-1}|\) and so \(|z^{-1}| = 1/|z|\). If \(z \in S^1\) it therefore follows that \(z^{-1} \in S^1\).

Let \(n \geq 1\) be an integer. The subset \(\mu_n\) of \(\mathbb{C}^\times\) consisting of all complex numbers \(x\) such that \(x^n = 1\) is a subgroup of \(\mathbb{C}^\times\), and in fact of \(S^1\), having \(n\) elements. It is called the \(n\)-th roots of unity. The proof is left as an exercise.

**Example 22.8.** Let \(F\) be a field and

\[
H = \left\{ \begin{pmatrix} 1 & a \\ 0 & 1 \end{pmatrix} : a \in F \right\}.
\]

Then \(H\) is a subgroup of \(GL_2(F)\).

**Definition 22.9.** Let \(G\) be a group. \(G\) is called cyclic if there is an element \(g \in G\) such that \(G = \{g^n : n \in \mathbb{Z}\}\); that is, any element of \(G\) is a power of \(g\). The element \(g\) is then called a generator of \(G\).

**Example 22.10.** Let \(G\) be any group. Let \(g \in G\) and define

\[
\langle g \rangle := \{g^n : n \in \mathbb{Z}\}.
\]

This is a cyclic subgroup of \(G\) (it may be finite or infinite).

**Example 22.11.** The group \(\mathbb{Z}\) is cyclic. As a generator we may take 1 (or \(-1\)).

**Example 22.12.** The group \((\mathbb{Z}/5\mathbb{Z})^\times = \{1, 2, 3, 4\}\) is cyclic. The elements 2, 3 are generators. The group \(\mathbb{Z}/8\mathbb{Z}^\times\) is not cyclic. One can check that the square of any element is 1.

### 23. The permutation and dihedral groups

23.1. **Permutation groups.**

**Definition 23.1.** A *permutation* of a set \(T\) is a bijective function \(f : T \to T\). We shall denote the set of permutations of \(T\) by \(S_T\). If \(T = \{1, 2, \cdots, n\}\) then we shall denote \(S_T\) as \(S_n\). It is called “the symmetric group on \(n\) letters”.

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The page is filled with mathematical content, mostly concerning group theory, with a focus on subgroups, cyclic groups, and permutation groups. It includes several examples and definitions that illustrate these concepts, such as the calculation of subgroups and the identification of cyclic subgroups within specific groups like \(\mathbb{C}^\times\) and \(\mathbb{Z}\). The text also delves into permutation groups, defining what constitutes a permutation and the symmetric group on a set of elements.
Proposition 23.2. For every non-empty set $T$, $S_T$ is a group under composition of functions. The cardinality of $S_n$ is $n!$.

Proof. The product of two permutations $f, g$ is their composition $f \circ g$; it is again a permutation. We have $[(f \circ g) \circ h](t) = (f \circ g)(h(g)) = f((g(h(t))) = f((g \circ h))(t)$. Thus, as functions, we have $(f \circ g) \circ h = f \circ (g \circ h)$ and so our operation is associative.

The identity is just the identity function. The inverse of a permutation $f$ is the inverse function $f^{-1}$, which satisfies $f \circ f^{-1} = f^{-1} \circ f = \text{Id}_T$.

Finally, to define a permutation $f$ on $\{1, 2, \ldots, n\}$ we can choose the image of 1 arbitrarily ($n$ choices), the image of 2 could be any element different from $f(1)$ ($n - 1$ choices), the image of 3 can be any elements different from the images of 1 and 2 ($n - 1$ choices), and so on. Altogether, we have $n \cdot (n - 1) \cdot (n - 1) \cdots 2 \cdot 1 = n!$ choices. \hfill $\Box$

Example 23.3. (1) For $n = 1$, $S_1$ consist of a single element and so is the trivial group.

(2) For $n = 2$ we have two permutations.
   (a) $\text{Id}$. $\text{Id}(1) = 1, \text{Id}(2) = 2$.
   (b) $\sigma$. $\sigma(1) = 2, \sigma(2) = 1$.

   We may also represent those permutations as tables:
   $$\text{Id} = \begin{pmatrix} 1 & 2 \\ 1 & 2 \end{pmatrix}, \quad \sigma = \begin{pmatrix} 1 & 2 \\ 2 & 1 \end{pmatrix}.$$

(3) For $n = 3$ we have 6 permutations. One of them is $\sigma$ given by $\sigma(1) = 2, \sigma(2) = 3, \sigma(3) = 1$, or in table form
   $$\sigma = \begin{pmatrix} 1 & 2 & 3 \\ 2 & 3 & 1 \end{pmatrix}.$$

   The table form is a better notation and we list all elements of $S_3$ in that form.

   $$\begin{pmatrix} 1 & 2 & 3 \\ 1 & 2 & 3 \end{pmatrix}, \quad \begin{pmatrix} 1 & 2 & 3 \\ 2 & 3 & 1 \end{pmatrix}, \quad \begin{pmatrix} 1 & 2 & 3 \\ 3 & 1 & 2 \end{pmatrix},$$

   $$\begin{pmatrix} 1 & 2 & 3 \\ 2 & 1 & 3 \end{pmatrix}, \quad \begin{pmatrix} 1 & 2 & 3 \\ 1 & 3 & 2 \end{pmatrix}, \quad \begin{pmatrix} 1 & 2 & 3 \\ 3 & 2 & 1 \end{pmatrix}.$$

   The first line is in fact a cyclic subgroup of $S_3$. It is the subgroup generated by $\begin{pmatrix} 1 & 2 & 3 \\ 2 & 3 & 1 \end{pmatrix}$

The groups $S_n$ are not commutative for $n \geq 3$. For example:

$$\begin{pmatrix} 1 & 2 & 3 \\ 2 & 1 & 3 \end{pmatrix} \begin{pmatrix} 1 & 2 & 3 \\ 2 & 3 & 1 \end{pmatrix} = \begin{pmatrix} 1 & 2 & 3 \\ 1 & 3 & 2 \end{pmatrix}, \quad \begin{pmatrix} 1 & 2 & 3 \\ 1 & 3 & 2 \end{pmatrix} \begin{pmatrix} 1 & 2 & 3 \\ 2 & 1 & 3 \end{pmatrix} = \begin{pmatrix} 1 & 2 & 3 \\ 3 & 2 & 1 \end{pmatrix}.$$

Here is another example of multiplication, in $S_5$ this time:

$$\begin{pmatrix} 1 & 2 & 3 & 4 & 5 \\ 3 & 2 & 1 & 4 & 5 \end{pmatrix} \begin{pmatrix} 1 & 2 & 3 & 4 & 5 \\ 5 & 2 & 4 & 3 & 1 \end{pmatrix} = \begin{pmatrix} 1 & 2 & 3 & 4 & 5 \\ 5 & 2 & 4 & 1 & 3 \end{pmatrix}.$$
23.2. **Cycles.** There is still more efficient notation for permutations in $S_n$. Fix $n \geq 1$. A cycle (in $S_n$) is an expression of the form

$$(a_1 \ a_2 \cdots a_t),$$

where $a_i \in \{1,2,\ldots,n\}$ are distinct elements. This expression is understood as the permutation $\sigma$ given by

$$\sigma(a) = \begin{cases} 
  a_{i+1} & a = a_i, i < n, \\
  a_1 & a = a_n, \\
  a & \text{else}.
\end{cases}$$

For example, the permutation $\left( \begin{array}{lll} 1 & 2 & 3 \\
 2 & 3 & 1 \end{array} \right)$ is the cycle $(1 2 3)$ and the permutation $\left( \begin{array}{lll} 1 & 2 & 3 \\
 2 & 1 & 3 \end{array} \right)$ is the cycle $(1 2)$. A cycle with two elements $(i \ j)$ (not necessarily consecutive) is called a transposition.

**Definition 23.4.** Let $G$ be a group. The order of $G$, denoted $|G|$, or $\#G$, is the number of elements of $G$ (written $\infty$ if not finite).

Let $g \in G$. The order of $g$ is defined as $|\langle g \rangle|$, the order of the cyclic group generated by $g$. It is also denoted by $o(g)$, or $\text{ord}(g)$.

**Lemma 23.5.** Let $g \in G$, $o(g)$ is the minimal positive integer $k$ such that $g^k = e$.

**Proof.** Let $k$ be the minimal integer such that $g^k = e$ ($\infty$ if such doesn’t exist).

Suppose first that $o(g)$ is finite, say equals $r$. Then the $r + 1$ elements $\{e, g, g^2, \ldots, g^r\}$ cannot be distinct and so $g^i = g^j$ for some $0 \leq i < j \leq r$. It follows that $g^{j-i} = e$ and so $k \leq j - i \leq r$. In particular $k$ is also finite. So $r$ is finite implies $k$ is finite and $k \leq r$.

Suppose now that $k$ is finite. Let $n$ be an integer and write $n = ak + b$ where $0 \leq b < k$. Then $g^n = (g^k)^a g^b = e^a g^b = g^b$. We conclude that $\langle g \rangle \subseteq \{e, g, \cdots, g^{k-1}\}$, and so $k$ is finite implies that $r$ is finite and $r \leq k$. \qed

**Example 23.6.** Let $(a_1 a_2 \cdots a_t)$ be a cycle. Its order is $t$.

Two cycles $\sigma, \tau$ are called disjoint if they contain no common elements. In this case, clearly $\sigma \tau = \tau \sigma$. Moreover, $(\sigma \tau)^n = \sigma^n \tau^n$ and since $\sigma^n$ and $\tau^n$ are disjoint, $(\sigma \tau)^n = Id$ if and only if $\sigma^n = Id$ and $\tau^n = Id$. Thus $o(\sigma) | n, o(\tau) | n$ and we deduce that $o(\sigma \tau)$ (namely, the least $n$ such that $(\sigma \tau)^n = Id$ is $\text{lcm}(o(\sigma), o(\tau))$. Arguing in the same way a little more generally we obtain:

**Lemma 23.7.** Let $\sigma_1, \ldots, \sigma_n$ be disjoint permutations of orders $r_1, \ldots, r_n$, respectively. Then the order of the permutation $\sigma_1 \circ \sigma_2 \circ \cdots \circ \sigma_n$ is $\text{lcm}(r_1, r_2, \ldots, r_n)$.

Combining this lemma with the following proposition allows us to calculate the order of every permutation very quickly.

**Proposition 23.8.** Every permutation is a product of disjoint cycles.
We shall not provide formal proof of this proposition, but illustrate it by examples.

**Example 23.9.** Consider the permutation \( \sigma = \begin{pmatrix} 1 & 2 & 3 & 4 & 5 \\ 3 & 4 & 1 & 5 & 2 \end{pmatrix} \). To write it as a product of cycles we begin by (1 and check where 1 goes to. It goes to 3. So we write (13 and check where 3 goes to. It goes to 1 and so we have (13). The first number we didn’t consider is 2. 2 goes to 4 and so we write (13)(24 and 4 goes to 5 and so we write (13)(245. Now, 5 goes to 2 and so we have \( \sigma = (13)(245 \). The order of \( \sigma \) is \( lcm(2, 3) = 6 \).

**Example 23.10.** Consider the permutation \( \sigma = \begin{pmatrix} 1 & 2 & 3 & 4 & 5 & 6 & 7 & 8 & 9 & 10 \\ 3 & 4 & 2 & 1 & 10 & 6 & 9 & 5 & 7 & 8 \end{pmatrix} \). This is written as a product of disjoint transposition as follows \( (1324)(79)(5108) \). To find this expression, we did the same procedure described above: We start with (1 , continue with (13 , because 1 goes to 3, and then with (132 , because 3 goes to 2. Then we find that 2 goes to 4 which goes to 1 and we have found (1324). The first number not in this list is 5 which goes to 10 and so we have (1324)(510. Since 10 goes to 8 and 8 to 5 we get now (1324)(5108). The first number not in this list is 6 that goes to 6 and that gives (1324)(5108)(6). We then continue with 7. Since 7 goes to 9 which goes to 7 we have (1324)(5108)(6)(79). We have considered all numbers and so \( \sigma = (1324)(5108)(6)(79) = (1324)(5108)(79) \). The order of \( \sigma \) is \( lcm(4, 3, 2) = 12 \).

**Example 23.11.** Suppose we want to find a permutation of order 10 in \( S_7 \). We simply take \( (12345)(67) \). If we want to find a permutation of order 10 in \( S_{10} \) we can take either \( (12345)(67) \) or \( (12345678910) \) (and all variants on this).

Finally we remark on the computation of \( \sigma^{-1} \) for a permutation \( \sigma \). If \( \sigma \) is given in the form of a table, for example:

\[
\sigma = \begin{pmatrix} 1 & 2 & 3 & 4 & 5 & 6 & 7 & 8 & 9 & 10 \\ 3 & 4 & 2 & 1 & 10 & 6 & 9 & 5 & 7 & 8 \end{pmatrix},
\]

then because \( \sigma(i) = j \iff \sigma^{-1}(j) = i \), the table describing \( \sigma^{-1} \) is the same table but read from the bottom to the top. That is

\[
\sigma^{-1} = \begin{pmatrix} 3 & 4 & 2 & 1 & 10 & 6 & 9 & 5 & 7 & 8 \\ 1 & 2 & 3 & 4 & 5 & 6 & 7 & 8 & 9 & 10 \end{pmatrix}.
\]

Only that we follow our convention and write the columns in the conventional order and so we get

\[
\sigma^{-1} = \begin{pmatrix} 1 & 2 & 3 & 4 & 5 & 6 & 7 & 8 & 9 & 10 \\ 4 & 3 & 1 & 2 & 8 & 6 & 9 & 10 & 7 & 5 \end{pmatrix}.
\]

If \( \sigma \) is a cycle, say \( \sigma = (i_1 i_2 \ldots i_{k-1} i_k) \), then \( \sigma^{-1} \) is easily seen to be \( (i_k i_{k-1} \ldots i_2 i_1) \). So,

\[
(i_1 i_2 \ldots i_{k-1} i_k)^{-1} = (i_k i_{k-1} \ldots i_2 i_1).
\]
Now, if $\sigma$ is a product of disjoint cycles, $\sigma = \sigma_1\sigma_2 \ldots \sigma_r$ then $\sigma^{-1} = \sigma_r^{-1} \ldots \sigma_2^{-1}\sigma_1^{-1}$ (by a generalization of the rule $(ab)^{-1} = b^{-1}a^{-1}$), but, since those cycles are disjoint and so commute, we can also write this as $\sigma^{-1} = \sigma_1^{-1}\sigma_2^{-1} \ldots \sigma_r^{-1}$. (This last manipulation is wrong if the cycles are not disjoint!) Thus, for example, the inverse of $\sigma = (1324)(5\ 10\ 8)(79)$ is $(4231)(8\ 10\ 5)(97)$, which we can also write, if we wish, as $(1423)(5\ 8\ 10)(79)$.

23.3. The Dihedral group. Consider a regular polygon with $n$ sides in the plane, $n \geq 3$, symmetric around $(0,0)$. The dihedral group $D_n$ is defined as the symmetries of the polygon. Let us number the vertices of the polygon by $1, 2, \ldots, n$ in the clockwise direction and say the first vertex 1 lies on the $x$-axis. One sees that every symmetry must permutes the vertices and in fact either maintains or reverses their order. In fact, if $\sigma$ is a symmetry then $\sigma(1) = j$ and $\sigma(2) = j + 1$ or $j - 1$ (where we understand $n + 1$ as 1 and 1 − 1 as $n$) and $\sigma$ is uniquely determined by these conditions. For example, the permutation $y$ given by the cycle $(1\ 2\ 3 \cdots n)$ is an element of the dihedral group rotating the polygon by angle $360^0/n$ in the clockwise direction (so if $t$ is a point on the boundary of the polygon such that the line from $(0,0)$ to $t$ forms an angle $\theta$ with the $x$-axis, then $y(t)$ is the point forming an angle $\theta - 360^0/n$). Another symmetry, $x$, is reflection through the $x$-axis. The symmetry $x$ is given as permutation by the product $(2\ n), (3\ n-1) \cdots (n/2 \ 2 + n/2)$ if $n$ is even and $(2\ n), (3\ n-1) \cdots ((n+1)/2 \ 1 + (n+1)/2)$ if $n$ is odd. In terms of angles, $x$ changes an angle $\theta$ to $-\theta$.

Theorem 23.12. The elements of the dihedral group are

$$D_n = \{e, y, \ldots, y^{n-1}, x, yx, y^2x, \ldots, y^{n-1}x\},$$

and the relations $x^2 = y^n = 1$ and $xyxy = 1$ hold. In particular, $D_n$ has $2n$ elements.

Proof. It is enough to show that any vertex $j \in \{1, 2, \ldots, n\}$ there is a unique element of the set $\{e, y, \ldots, y^{n-1}, yx, y^2x, \ldots, y^{n-1}x\}$ that takes 1 to $j$ and 2 to $j + 1$ and there is a unique element taking 1 to $j$ and 2 to $j - 1$. This shows both that every element of $D_n$ is in the list and that all elements of the list are different.

We calculate that

$$y^a(1) = a + 1, \quad y^a(2) = a + 2,$$

and

$$y^a x(1) = a + 1, \quad y^a x(2) = y^a(n) = a.$$

This proves our claims.

The relations $x^2 = y^n = 1$ are evident. We check that $xyxy = 1$, by checking that $xyxy(j) = j$ for $j = 1, 2$. We have $xyxy(1) = xyx(2) = xy(n) = x(1) = 1$ and $xyxy(2) = xyx(3) = xy(n-1) = x(n) = 2$. □

The nature of the symmetries $1, y, \ldots, y^{n-1}$ is clear: $y^j$ rotates clockwise by angle $j \cdot 360^0/n$. 
Theorem 24.2. Let \( 0 \leq j < n \). The element \( y^jx \) is a reflection through the line forming an angle \(-j \cdot 360^0/2n\) with the x-axis.

Proof. The symmetry \( y^jx \) is not trivial. If it fixes an angle \( \theta \) it must be reflection through the line with that angle. Note that \( y^jx \) sends the angle \( \theta \) to \(-\theta\) and then adds \(-j \cdot 360^0/\) so the equation is \( \theta = -\theta - j \cdot 360^0/\) (mod 360). That is \( \theta = -j \cdot 360^0/2n \).

24. The theorem of Lagrange

24.1. Cosets. Let \( H < G \) be a subgroup of \( G \). A left coset of \( H \) in \( G \) is a subset of the form

\[
gH := \{gh : h \in H\},
\]

for some \( g \in G \). The set \( gH \) is called the left coset of \( g \); \( g \) is called a representative of the coset \( gH \).

Example 24.1. Consider the subgroup \( H \) of \( S_3 \) given by \( \{1, (123), (132)\} \). Here are some cosets: \( H = 1H = (123)H = (132)H, \) \( (12)H = (13)H = (23)H = \{(12), (23), (13)\} \). We leave the verification to the reader.

Lemma 24.2. Let \( H \) be a subgroup of \( G \).

1. Two left cosets are either equal or disjoint.
2. Let \( g_1H, g_2H \) be two left cosets. The following are equivalent: (i) \( g_1H = g_2H \); (ii) \( g_1 \in g_2H \); (iii) \( g_2^{-1}g_1 \in H \).

Proof. Suppose that \( g_1H \cap g_2H \neq \emptyset \), so for some \( h_1, h_2 \) we have \( g_1h_1 = g_2h_2 \). We prove that \( g_1H \subseteq g_2H \). By symmetry we also have \( g_2H \subseteq g_1H \) and so \( g_1H = g_2H \).

Let \( h \in H \). Then \( g_1h = ((g_2h_2)h_1^{-1})h = g_2(h_2h_1^{-1}h) \in g_2H \).

We now prove the equivalence of the assertions (i) - (iii). Suppose (i) holds. Then \( g_1 = g_1e \in g_2H \) and (ii) holds. Suppose (ii) holds; say \( g_1 = g_2h \). Then \( g_2^{-1}g_1 = h \in H \) and (iii) holds. Suppose that (iii) holds; \( g_2^{-1}g_1 = h \) for some \( h \in H \). Then \( g_1 = g_2h \) and so \( g_1H \cap g_2H \neq \emptyset \). By what we have proved in the first part, \( g_1H = g_2H \).

Remark 24.3. The Lemma and its proof should be compared with Lemma 19.3. In fact, since \( R \) is an abelian group and an ideal \( I \) is a subgroup, that lemma is special case of the lemma above.

Corollary 24.4. \( G \) is a disjoint union of cosets of \( H \). Let \( \{g_i : i \in I\} \) be a set of elements of \( G \) such that each coset has the form \( g_iH \) for a unique \( g_i \). That is, \( G = \bigsqcup_{i \in I} g_iH \). Then the \( \{g_i : i \in I\} \) are called a complete set of representatives.

In the same manner one defines a right coset of \( H \) in \( G \) to be a subset of the form \( Hg = \{hg : h \in H\} \) and Lemma 24.2 holds for right cosets with the obvious modifications. Two right cosets
are either equal or disjoint and the following are equivalent: (i) $H g_1 = H g_2$; (ii) $g_1 \in H g_2$; (iii) $g_1 g_2^{-1} \in H$. Thus, the Corollary holds true for right cosets as well.

We remark that the intersection of a left coset and a right coset may be non-empty, yet not a coset itself. For example, take $H = \{1, (12)\}$ in $S_3$. We have the following table.

<table>
<thead>
<tr>
<th>$g$</th>
<th>$gH$</th>
<th>$Hg$</th>
</tr>
</thead>
<tbody>
<tr>
<td>1</td>
<td>${1, (12)}$</td>
<td>${1, (12)}$</td>
</tr>
<tr>
<td>(12)</td>
<td>${(12), 1}$</td>
<td>${(12), 1}$</td>
</tr>
<tr>
<td>(13)</td>
<td>${(13), (123)}$</td>
<td>${(13), (132)}$</td>
</tr>
<tr>
<td>(23)</td>
<td>${(23), (132)}$</td>
<td>${(23), (123)}$</td>
</tr>
<tr>
<td>(123)</td>
<td>${(123), (13)}$</td>
<td>${(123), (23)}$</td>
</tr>
<tr>
<td>(132)</td>
<td>${(132), (23)}$</td>
<td>${(132), (13)}$</td>
</tr>
</tbody>
</table>

The table demonstrates that indeed any two left (resp. right) cosets are either equal or disjoint, but the intersection of a left coset with a right coset may be non-empty and properly contained in both.

24.2. Lagrange’s theorem.

**Theorem 24.5.** Let $G$ be a finite group and $H$ a subgroup of $G$. Then,

$$|H| \mid |G|.$$ 

Moreover, let $\{g_i : i \in I\}$ be a complete set of representatives for the cosets of $H$, then $|I| = \frac{|G|}{|H|}$. In particular, the cardinality of $I$ does not depend on the choice of a complete set of representatives. It is called the index of $H$ in $G$.

**Proof.** We have,

$$G = \coprod_{i \in I} g_i H.$$ 

Let $a, b \in G$. We claim that the function

$$f : aH \to bH, \quad x \mapsto ba^{-1}x,$$

is a well defined bijection. First, $x = ah$ for some $h$ and so $ba^{-1}x = bh \in bH$ and so the map is well defined. It is surjective, because given an element $y \in bH$, say $y = bh$ it is the image of $ah$. The map is also injective: if $ba^{-1}x_1 = ba^{-1}x_2$ then multiplying both sides by $ab^{-1}$ we get $x_1 = x_2$.

We conclude that each coset $g_i H$ has the same number of elements, which is exactly the number of elements in $H = eH$. We get therefore that

$$|G| = |H| \cdot |I|.$$ 

That completes the proof.

Here are some applications of Lagrange’s theorem:
Let $G$ be a finite group of prime order $p$. Then $G$ is cyclic; in fact, every element of $G$ that is not the identity generates $G$.

Indeed, Let $g \neq e$. Then $H = \langle g \rangle$ is a non-trivial subgroup. So $|H| > 1$ and divides $p$. It follows that $|H| = |G|$ and so that $\langle g \rangle = G$.

(2) In a similar vein, we conclude that a group of order 6 say, cannot have elements of order 4, or 5, or any order not dividing 6. This follows immediately from Lagrange’s theorem, keeping in mind that $\text{ord}(g) = |\langle g \rangle|$.

25. HOMOMORPHISMS AND ISOMORPHISMS

25.1. HOMOMORPHISMS OF GROUPS.

Definition 25.1. Let $G, H$ be groups and

$$f : G \to H,$$

a function. The function $f$ is called a group homomorphism, if

$$f(g_1g_2) = f(g_1)f(g_2), \quad \forall g_1, g_2 \in G.$$ 

In that case, we define the kernel of $f$ as:

$$\text{Ker}(f) = \{g \in G : f(g) = e_H\}.$$ 

Lemma 25.2. Let $f : G \to H$ be a group homomorphism. Then:

1. $f(e_G) = e_H$;
2. $f(g^{-1}) = f(g)^{-1}$;
3. The image of $f$ is a subgroup of $H$.

Proof. We have $f(e_G) = f(e_G e_G) = f(e_G)f(e_G)$. Multiplying (in $H$) both sides by $f(e_G)^{-1}$ we find $e_H = f(e_G)$. Now, $e_H = f(e_G) = f(gg^{-1}) = f(g)f(g^{-1})$, which shows that $f(g^{-1}) = f(g)^{-1}$.

Finally, we show that $\text{Im}(f)$ is a subgroup of $H$. Note that $e_H = f(e_G) \in \text{Im}(f)$. If $h_1, h_2 \in \text{Im}(f)$, say $h_i = f(g_i)$ then $h_1 h_2 = f(g_1 g_2)$ and $h_1^{-1} = f(g_1^{-1})$. This shows that $h_1 h_2, h_1^{-1} \in \text{Im}(f)$. \qed

Proposition 25.3. Let $f : G \to H$ be a group homomorphism. $\text{Ker}(f)$ is a subgroup of $G$. The homomorphism $f$ is injective if and only if $\text{Ker}(f) = \{e_G\}$.

Proof. First, we proved that $f(e_G) = e_H$ and so $e_G \in \text{Ker}(f)$. Next, if $g_1, g_2 \in \text{Ker}(f)$ then $f(g_1) = f(g_2) = e_H$ and so $f(g_1 g_2) = f(g_1)f(g_2) = e_He_H = e_H$. Therefore, $g_1 g_2 \in H$.

Finally, if $g \in \text{Ker}(f)$ then $f(g^{-1}) = f(g)^{-1} = e_H^{-1} = e_H$ and so $g^{-1} \in \text{Ker}(f)$ as well.

Suppose $f$ is injective. Then, since $f(e_g) = e_H$, $e_G$ is the only element mapping to $e_H$ and so $\text{Ker}(f) = \{e_G\}$. Conversely, suppose $\text{Ker}(f) = \{e_G\}$ and $f(g_1) = f(g_2)$. Then $e_H = f(g_1)^{-1}f(g_2) = f(g_1^{-1})f(g_2) = f(g_1^{-1}g_2)$. That means that $g_1^{-1}g_2 \in \text{Ker}(f)$ and so $g_1^{-1}g_2 = e_G$. That is, $g_1 = g_2$. \qed
25.2. Isomorphism.

Definition 25.4. A group homomorphism \( f : G \to H \) is called an isomorphism if it is bijective.

As in the case of rings, one verifies that if \( f \) is an isomorphism, the inverse function \( g = f^{-1} \) is automatically a homomorphism and so an isomorphism as well. Also, one easily checks that a composition of group homomorphisms is a group homomorphism. It follows that being isomorphic is an equivalence relation on groups. Cf. §20.1.

Example 25.5. Let \( n \) be a positive integer. Any two cyclic groups of order \( n \) are isomorphic.

Indeed, suppose that \( G = \langle g \rangle, H = \langle h \rangle \) are cyclic groups of order \( n \). Define, for any integer \( a \), \( f(g^a) = h^a \).

This is well defined; if \( g^a = g^b \) then \( g^{a-b} = e_G \) and so \( n | (a - b) \). Thus, \( a = b + kn \) and \( f(g^a) = h^a = h^b(h^n)^k = h^b = f(g^b) \). Obviously \( f \) is a surjective homomorphism; \( f \) is also injective, because \( f(g^a) = h^a = e_H \) implies that \( n | a \) and so \( g^a = e_G \).

In particular, we conclude that any cyclic group of order \( n \) is isomorphic to the group \( \mathbb{Z}/n\mathbb{Z} \) with addition.

Example 25.6. Let \( p \) be a prime number then any two groups of order \( p \) are isomorphic. Indeed, we’ve seen that such groups are necessarily cyclic.

Theorem 25.7. (Cayley) Let \( G \) be a finite group of order \( n \) then \( G \) is isomorphic to a subgroup of \( S_n \).

Proof. Let \( g \in G \) and let \( \sigma_g : G \to G, \sigma_g(a) = ga \).

We claim that \( \sigma_g \) is a permutation. It is injective, because \( \sigma_g(a) = \sigma_g(b) \Rightarrow ga = gb \Rightarrow a = b \). It is surjective, because for any \( b \in G \), \( \sigma_g(g^{-1}b) = b \).

Identifying the permutations of \( G \) with \( S_n \) (just call the elements of \( G, 1, 2, 3, \ldots \)), we got a map \( G \to S_n, g \mapsto \sigma_g \).

This map is a homomorphism of groups: \( \sigma_{gh}(a) = gha = \sigma_g(\sigma_h(a)) \). That is, \( \sigma_{gh} = \sigma_g \circ \sigma_h \).

This homomorphism is injective: if \( \sigma_g \) is the identity permutation then \( \sigma_g(e) = e \) and that implies \( ge = e \), that is \( g = e \). We get that \( G \) is isomorphic to its image, which is a subgroup of \( S_n \), under this homomorphism. \( \square \)

Remark 25.8. We were somewhat informal about identifying the permutations of \( G \) with \( S_n \). A more rigorous approach is the following.

Lemma 25.9. Let \( T, Z \) be sets and \( f : T \to Z \) a bijection. The group of permutations of \( T \) and \( Z \) are isomorphic.
Proof. Let \( \sigma \in S_T \), a permutation of \( T \). Then \( f \circ \sigma \circ f^{-1} \) is a function from \( Z \) to itself, and being a composition of bijections is a bijection itself. We shall write more simply \( f\sigma f^{-1} \) for \( f \circ \sigma \circ f^{-1} \). We therefore got a function

\[
S_T \to S_Z, \quad \sigma \mapsto f\sigma f^{-1}.
\]

We claim that \( \sigma \mapsto f\sigma f^{-1} \) is a homomorphism. Indeed, given \( \sigma_1, \sigma_2 \in S_T \) we have \( f\sigma_1 f^{-1} = (f\sigma_1 f^{-1})(f\sigma_2 f^{-1}) \). Moreover, it is easy to write an inverse to this homomorphism,

\[
S_Z \to S_T, \quad \tau \mapsto f^{-1}\tau f.
\]

Therefore, we found a bijective homomorphism \( S_T \to S_Z \), which shows those two permutation groups are isomorphic. \( \square \)

26. Group actions on sets

26.1. Basic definitions. Let \( G \) be a group and let \( S \) be a non-empty set. We say that \( G \) acts on \( S \) if we are given a function

\[
G \times S \to S, \quad (g, s) \mapsto g \star s,
\]

such that;

(i) \( e \star s = s \) for all \( s \in S \);

(ii) \( (g_1g_2) \star s = g_1 \star (g_2 \star s) \) for all \( g_1, g_2 \in G \) and \( s \in S \).

Given an action of \( G \) on \( S \) we can define the following sets. Let \( s \in S \). Define the Orbit of \( s \)

\[
\text{Orb}(s) = \{ g \star s : g \in G \}.
\]

Note that \( \text{Orb}(s) \) is a subset of \( S \), equal to all the images of the element \( s \) under the action of the elements of the group \( G \). We also define the stabilizer of \( s \) to be

\[
\text{Stab}(s) = \{ g \in G : g \star s = s \}.
\]

Note that \( \text{Stab}(s) \) is a subset of \( G \). In fact, it is a subgroup, as Lemma 26.1 states.

26.2. Basic properties.

Lemma 26.1. (1) Let \( s_1, s_2 \in S \). We say that \( s_1 \) is related to \( s_2 \), i.e., \( s_1 \sim s_2 \), if there exists \( g \in G \) such that \( g \star s_1 = s_2 \). This is an equivalence relation. The equivalence class of \( s_1 \) is its orbit \( \text{Orb}(s_1) \).

(2) Let \( s \in S \). The set \( \text{Stab}(s) \) is a subgroup of \( G \).

(3) Suppose that both \( G \) and \( S \) have finitely many elements. Then

\[
|\text{Orb}(s)| = \frac{|G|}{|\text{Stab}(s)|}.
\]
Let \( e \star s = s \) and hence \( s \sim s \), meaning the relation is reflexive. Second, if \( s_1 \sim s_2 \) then for a suitable \( g \in G \) we have \( g \star s_1 = s_2 \). Therefore, \( g^{-1} \star (g \star s_1) = g^{-1} \star s_2 \) and \( (g^{-1}) \star s_1 = g^{-1} \star s_2 \). It follows that, \( e \star s_1 = g^{-1} \star s_2 \) and so, \( s_1 = g^{-1} \star s_2 \), which implies that \( s_2 \sim s_1 \).

It remains to show the relation is transitive. If \( s_1 \sim s_2 \) and \( s_2 \sim s_3 \) then for suitable \( g_1, g_2 \in G \) we have \( g_1 \star s_1 = s_2 \) and \( g_2 \star s_2 = s_3 \). Therefore, \( (g_2 g_1) \star s_1 = g_2 \star (g_1 \star s_1) = g_2 \star s_2 = s_3 \), and hence \( s_1 \sim s_3 \).

Moreover, by the very definition, the equivalence class of an element \( s_1 \) of \( S \) is all the elements of the form \( g \star s_1 \) for some \( g \in G \), namely, \( \text{Orb}(s_1) \).

(2) Let \( H = \text{Stab}(s) \). We have to show that: (i) \( e \in H \), (2) if \( g_1, g_2 \in H \) then \( g_1 g_2 \in H \), and (iii) if \( g \in H \) then \( g^{-1} \in H \).

First, by definition of group action, we have \( e \star s = s \). Therefore, \( e \in H \). Next, suppose that \( g_1, g_2 \in H \), i.e., \( g_1 \star s = s \) and \( g_2 \star s = s \). Then, \( (g_1 g_2) \star s = g_1 \star (g_2 \star s) = g_1 \star s = s \), which proves that \( g_1 g_2 \in H \). Finally, if \( g \in H \) then \( g \star s = s \) and so \( g^{-1} \star (g \star s) = g^{-1} \star s \). That is, \( (g^{-1}g) \star s = g^{-1} \star s \) and so \( e \star s = g^{-1} \star s \), or \( s = g^{-1} \star s \), and therefore \( g^{-1} \in H \).

(3) We claim that there exists a bijection between the left cosets of \( H \) and the orbit of \( s \). If we show that, then by Lagrange’s theorem,

\[
|\text{Orb}(s)| = \text{no. of left cosets of } H = \text{index of } H = |G|/|H|.
\]

Define a function

\[
\{\text{left cosets of } H\} \xrightarrow{\phi} \text{Orb}(s),
\]

by

\[
\phi(gH) = g \star s.
\]

We claim that \( \phi \) is a well-defined bijection. First

**Well-defined:** Suppose that \( g_1 H = g_2 H \). We need to show that the rule \( \phi \) would give the same result whether we take the representative \( g_1 \) or the representative \( g_2 \) to the coset, that is, we need to show \( g_1 \star s = g_2 \star s \). Note that \( g_1^{-1} g_2 \in H \), i.e., \( (g_1^{-1} g_2) \star s = s \). We get \( g_1 \star s = g_1 \star ((g_1^{-1} g_2) \star s) = (g_1 (g_1^{-1} g_2)) \star s = g_2 \star s \).

**\( \phi \) is surjective:** Let \( t \in \text{Orb}(s) \) then \( t = g \star s \) for some \( g \in G \). Thus, \( \phi(gH) = g \star s = t \), and we get that \( \phi \) is surjective.

**\( \phi \) is injective:** Suppose that \( \phi(g_1 H) = \phi(g_2 H) \). We need to show that \( g_1 H = g_2 H \). Indeed, \( \phi(g_1 H) = \phi(g_2 H) \) implies \( g_1 \star s = g_2 \star s \) and so that \( g_2^{-1} \star (g_1 \star s) = g_2^{-1} \star (g_2 \star s) \); that is, \( (g_2^{-1} g_1) \star s = (g_2^{-1} g_2) \star s \) and so \( (g_2^{-1} g_1) \star s = e \star s = s \). Therefore, \( g_2^{-1} g_1 \in \text{Stab}(s) = H \) and hence \( g_1 H = g_2 H \). \( \square \)

**Corollary 26.2.** The set \( S \) is a disjoint union of orbits.

**Proof.** The orbits are the equivalence classes of the equivalence relation \( \sim \) defined in Lemma 26.1. Any equivalence relation partitions the set into disjoint equivalence classes. \( \square \)
26.3. Some examples.

Example 26.3. Let $G$ be the group of real numbers $\mathbb{R}$. The group operation is addition. Let $S$ be the set of points on the sphere in $\mathbb{R}^3$ of radius 1 about the origin. The group $\mathbb{R}$ acts by rotating around the $z$-axis. An element $r \in \mathbb{R}$ rotates by degree $r$. For every point $s \in S$, different from the poles, the stabilizer is $2\pi\mathbb{Z}$. For the poles the stabilizer is $\mathbb{R}$. The orbit of every point is the altitude line on which it lies.

Example 26.4. Let $G$ be a group and $H$ a subgroup of $G$. Then $H$ acts on $G$ by

$$H \times G \to G, \quad (h, g) \mapsto hg.$$ 

Here $H$ plays the role of the group and $G$ the role of the set in the definition. This is indeed a group action: $e_H g = g$ for all $g \in G$, because by definition $e_H = e_G$. Also, $h_1(h_2)g = (h_1h_2)g$ is nothing but the associative law.

The orbit of $g \in G$ is

$$\text{Orb}(g) = \{hg : h \in H\} = Hg.$$ 

That is, the orbits are the right cosets of $H$. We have that $G$ is a disjoint union of orbits, namely, a disjoint union of cosets. The stabilizer of any element $g \in G$ is $\{e\}$. The formula we have proven, $|\text{Orb}(g)| = |H|/|\text{Stab}(g)|$, gives us $|Hg| = |H|$ for any $g \in G$, and we see that we have another point of view on Lagrange's theorem.

Example 26.5. We consider a roulette with $n$ sectors and write $n = i_1 + \cdots + i_k$, for some positive (and fixed) integers $i_1, \ldots, i_k$. We suppose we have different colors $c_1, \ldots, c_k$ and we color $i_1$ sectors of the roulette by the color $c_1$, $i_2$ sectors by the color $c_2$ and so on. The sectors can be chosen as we wish and so there are many possibilities. We get a set $S$ of colored roulettes. Now, we turn the roulette $a$ steps clockwise, say, then we get another colored roulette, usually with different coloring. Nonetheless, it is natural to view the two coloring as the same, since “they only depend on your point of view”. We may formalize this by saying that the group $\mathbb{Z}/n\mathbb{Z}$ acts on $S$, $a$ acts on a colored roulette by turning it $a$ steps clockwise, and by saying that we are interested in the number of orbits for this action.

Example 26.6. Let $G$ be the dihedral group $D_8$. Recall that $G$ is the group of symmetries of a regular octagon in the plane.

$$G = \{e, y, y^2, \ldots, y^7, x, yx, y^2x, \ldots, y^7x\},$$ 

where $y$ is rotation clockwise by angle $2\pi/8$ and $x$ is reflection through the $x$-axis. We have the relations

$$x^2 = y^8 = e, \quad xyxy = 1.$$ 

We let $S$ be the set of colorings of the octagon ($= \text{necklaces laid on the table}$) having 4 red vertices (rubies) and 4 green vertices (sapphires). The group $G$ acts on $S$ by its action on the octagon.
For example, the coloring $s_0$ consisting of alternating green and red is certainly preserved under $x$ and under $y^2$. Therefore, the stabilizer of $s_0$ contains at least the set of eight elements

\[(26.1)\quad \{e, y^2, y^4, y^6, x, y^2x, y^4x, y^6x\}.\]

Remember that the stabilizer is a subgroup and, by Lagrange’s theorem, of order dividing $16 = |G|$. On the other hand, $\text{Stab}(s_0) \neq G$ because $y \notin \text{Stab}(s_0)$. It follows that the stabilizer has exactly 8 elements and is equal to the set in (26.1).

Let $H$ be the stabilizer of $s_0$. According to Lemma 26.1 the orbit of $s_0$ is in bijection with the left cosets of $H = \{e, y^2, y^4, y^6, x, y^2x, y^4x, y^6x\}$. By Lagrange’s theorem there are two cosets. For example, $H$ and $gH$ are distinct cosets. The proof of Lemma 26.1 tells us how to find the orbit: it is the set $\{s_0, gs_0\}$, which is of course quite clear if you think about it.

27. **The Cauchy-Frobenius Formula**

**Theorem 27.1. (CFF)**\(^{11}\) Let $G$ be a finite group acting on a finite set $S$. Let $N$ be the number of orbits of $G$ in $S$. Define

\[I(g) = |\{s \in S : g \ast s = s\}|\]

(the number of elements of $S$ fixed by the action of $g$). Then \(^{12}\)

\[N = \frac{1}{|G|} \sum_{g \in G} I(g).\]

**Remark 27.2.** Note that $I(g)$ is the number of fixed points for the action of $g$ on $S$. Thus, the CFF can be interpreted as saying that the number of orbits is the average number of fixed points (though this does not make the assertion more obvious).

**Proof.** We define a function

\[T : G \times S \rightarrow \{0, 1\}, \quad T(g, s) = \begin{cases} 1 & g \ast s = s \\ 0 & g \ast s \neq s. \end{cases}\]

Note that for a fixed $g \in G$ we have

\[I(g) = \sum_{s \in S} T(g, s),\]

\(^{11}\)This is also sometimes called Burnside’s formula.

\(^{12}\)The sum appearing in the formula means just that: If you write $G = \{g_1, \ldots, g_n\}$ then $\sum_{g \in G} I(g)$ is $\sum_{i=1}^n I(g_i) = I(g_1) + I(g_2) + \cdots + I(g_n)$. The double summation $\sum_{g \in G} \sum_{s \in S} T(g, s)$ appearing in the proof means that if we write $S = \{s_1, \ldots, s_m\}$ then the double sum is $T(g_1, s_1) + T(g_1, s_2) + \cdots + T(g_1, s_m) + T(g_2, s_1) + T(g_2, s_2) + \cdots + T(g_2, s_m) + \cdots + T(g_n, s_1) + T(g_n, s_2) + \cdots + T(g_n, s_m)$.
and that for a fixed \( s \in S \) we have
\[
|\text{Stab}(s)| = \sum_{g \in G} T(g, s).
\]

Let us fix representatives \( s_1, \ldots, s_N \) for the \( N \) disjoint orbits of \( G \) in \( S \). Now,
\[
\sum_{g \in G} I(g) = \sum_{g \in G} \left( \sum_{s \in S} T(g, s) \right) = \sum_{s \in S} \left( \sum_{g \in G} T(g, s) \right)
\]
\[
= \sum_{s \in S} |\text{Stab}(s)| = \sum_{s \in S} \frac{|G|}{|\text{Orb}(s)|}
\]
\[
= \sum_{i=1}^{N} \frac{|G|}{|\text{Orb}(s_i)|} \cdot |\text{Orb}(s_i)| = \sum_{i=1}^{N} |G|
\]
\[
= N \cdot |G|.
\]

\[ \square \]

Remark 27.3. If \( N \), the number of orbits, is equal 1 we say that \( G \) acts transitively on \( S \). It means exactly that: For every \( s_1, s_2 \in S \) there exists \( g \in G \) such that \( g \star s_1 = s_2 \). Note that if \( G \) and \( S \) are finite then if \( G \) acts transitively then the number of elements in \( S \) divides the number of elements in \( G \),
\[
|S| \mid |G|,
\]
because, if \( S = \text{Orb}(s) \) then \( |S| = |G|/|\text{Stab}(s)| \).

Corollary 27.4. Let \( G \) be a finite group acting transitively on a finite set \( S \). Suppose that \( |S| > 1 \). Then there exists \( g \in G \) without fixed points.

Proof. By contradiction. Suppose that every \( g \in G \) has a fixed point in \( S \). That is, suppose that for every \( g \in G \) we have
\[
I(g) \geq 1.
\]
Since \( I(e) = |S| > 1 \) we have that
\[
\sum_{g \in G} I(g) > |G|.
\]
By Cauchy-Frobenius formula, the number of orbits \( N \) is greater than 1. Contradiction. \[ \square \]
27.1. Some applications to Combinatorics.

Example 27.5. How many roulettes with 11 wedges painted 2 blue, 2 green and 7 red are there when we allow rotations?

Let $S$ be the set of painted roulettes. Let us enumerate the sectors of a roulette by the numbers $1, \ldots, 11$. The set $S$ is a set of $\binom{11}{2} \binom{9}{2} = 1980$ elements (choose which 2 are blue, and then choose out of the nine left which 2 are green).

Let $G$ be the group $\mathbb{Z}/11\mathbb{Z}$. It acts on $S$ by rotations. The element 1 rotates a painted roulette by angle $2\pi/11$ anti-clockwise. The element $n$ rotates a painted roulette by angle $2n\pi/11$ anti-clockwise. We are interested in $N$ – the number of orbits for this action. We use CFF.

The identity element always fixes the whole set. Thus $I(0) = 1980$. We claim that if $1 \leq i \leq 10$ then $i$ doesn’t fix any element of $S$. We use the following fact that we have proved before: Let $G$ be a finite group of prime order $p$. Let $g \neq e$ be an element of $G$. Then $\langle g \rangle = G$.

Suppose that $1 \leq i \leq 10$ and $i$ fixes $s$. Then so does $(i) = \mathbb{Z}/11\mathbb{Z}$ (the stabilizer is a subgroup). But any coloring fixed under rotation by 1 must be single colored! Contradiction.

Applying CFF we get

$$N = \frac{1}{11} \sum_{n=0}^{10} I(n) = \frac{1}{11} \cdot 1980 = 180.$$ 

Example 27.6. How many roulettes with 12 wedges painted 2 blue, 2 green and 8 red are there when we allow rotations?

Let $S$ be the set of painted roulettes. Let us enumerate the sectors of a roulette by the numbers $1, \ldots, 12$. The set $S$ is a set of $\binom{12}{2} \binom{10}{2} = 2970$ elements (choose which 2 are blue, and then choose out of the ten left which 2 are green).

Let $G$ be the group $\mathbb{Z}/12\mathbb{Z}$. It acts on $S$ by rotations. The element 1 rotates a painted roulette by angle $2\pi/12$ anti-clockwise. The element $n$ rotates a painted roulette by angle $2n\pi/12$ anti-clockwise. We are interested in $N$ – the number of orbits for this action. We use CFF.

The identity element always fixes the whole set. Thus $I(0) = 2970$. We claim that if $1 \leq i \leq 11$ and $i \neq 6$ then $i$ doesn’t fix any element of $S$. Indeed, suppose that $i$ fixes a painted roulette. Say in that roulette the $r$-th sector is blue. Then so must be the $i+r$ sector (because the $r$-th sector goes under the action of $i$ to the $r+i$-th sector). Therefore so must be the $r+2i$ sector. But there are only 2 blue sectors! The only possibility is that the $r+2i$ sector is the same as the $r$ sector, namely, $i = 6$.

If $i$ is equal to 6 and we enumerate the sectors of a roulette by the numbers $1, \ldots, 12$ we may write $i$ as the permutation

$$(1\ 7)(2\ 8)(3\ 9)(4\ 10)(5\ 11)(6\ 12).$$
In any coloring fixed by \( i = 6 \) the colors of the pairs \((1, 7), (2, 8), (3, 9), (4, 10), (5, 11)\) and \((6, 12)\) must be the same. We may choose one pair for blue, one pair for green. The rest would be red. Thus there are \(30 = 6 \cdot 5\) possible choices. We summarize:

<table>
<thead>
<tr>
<th>element ( g )</th>
<th>( I(g) )</th>
</tr>
</thead>
<tbody>
<tr>
<td>0</td>
<td>2970</td>
</tr>
<tr>
<td>( i \neq 6 )</td>
<td>0</td>
</tr>
<tr>
<td>( i = 6 )</td>
<td>30</td>
</tr>
</tbody>
</table>

Applying CFF we get that there are

\[ N = \frac{1}{12}(2970 + 30) = 250 \]

different roulettes.

**Example 27.7.** In this example \( S \) is the set of necklaces made of four rubies and four sapphires laid on the table. We ask how many necklaces there are when we allow rotations and flipping-over.

We may talk of \( S \) as the colorings of a regular octagon, four vertices are green and four are red. The group \( G = D_8 \) acts on \( S \) and we are interested in the number of orbits for the group \( G \).

The results are the following

<table>
<thead>
<tr>
<th>element ( g )</th>
<th>( I(g) )</th>
</tr>
</thead>
<tbody>
<tr>
<td>( e )</td>
<td>70</td>
</tr>
<tr>
<td>( y, y^4, y^5, y^6 )</td>
<td>0</td>
</tr>
<tr>
<td>( y^2, y^6 )</td>
<td>2</td>
</tr>
<tr>
<td>( y^4 )</td>
<td>6</td>
</tr>
<tr>
<td>( xy^i ) for ( i = 0, \ldots, 7 )</td>
<td>6</td>
</tr>
</tbody>
</table>

We explain how the entries in the table are obtained:

The identity always fixes the whole set \( S \). The number of elements in \( S \) is \( \binom{8}{4} = 70 \) (choosing which 4 would be green).

The element \( y \) cannot fix any coloring, because any coloring fixed by \( y \) must have all sections of the same color (because \( y = (1 \ 2 \ 3 \ 4 \ 5 \ 6 \ 7 \ 8) \)). If \( y^r \) fixes a coloring \( s_0 \) so does \( (y^r)^r = y^{r^2} \) because the stabilizer is a subgroup. Apply that for \( r = 3, 5, 7 \) to see that if \( y^r \) fixes a coloring so does \( y^r \), which is impossible.  

Now, \( y^2 \), written as a permutation, is \((1 \ 3 \ 5 \ 7)(2 \ 4 \ 6 \ 8)\). We see that if, say 1 is green so are 3, 5, 7 and the rest must be red. That is, all the freedom we have is to choose whether the cycle \((1 \ 3 \ 5 \ 7)\) is green or red. This gives us two colorings fixed by \( y^2 \). The same rational applies to \( y^6 = (8 \ 6 \ 4 \ 2)(7 \ 5 \ 3 \ 1) \).

\[ y^{32} = y^8 = e \] because \( y^8 = e \), etc.
Consider now $y^4$. It may be written in permutation notation as $(1 5)(2 6)(3 7)(4 8)$. In any coloring fixed by $y^4$ each of the cycles $(1 5)(2 6)(3 7)$ and $(4 8)$ must be single colored. There are thus $\binom{4}{2} = 6$ possibilities (Choosing which 2 out of the four cycles would be green).

It remains to deal with the elements $x y^i$. We recall that these are all reflections. There are two kinds of reflections. One may be written using permutation notation as
\[(i_1 i_2)(i_3 i_4)(i_5 i_6)\]
(with the other two vertices being fixed. For example $x = (2 8)(3 7)(4 6)$ is of this form). The other kind is of the form
\[(i_1 i_2)(i_3 i_4)(i_5 i_6)(i_7 i_8).\]
(For example $x y = (1 8)(2 7)(3 6)(4 5)$ is of this sort). Whatever is the case, one uses similar reasoning to deduce that there are 6 colorings preserved by a reflection.

One needs only apply CFF to get that there are
\[N = \frac{1}{16}(70 + 2 \cdot 2 + 6 + 8 \cdot 6) = 8\]
distinct necklaces.

**Example 27.8.** Consider a tetrahedron with faces marked 1, 2, 3, 4. It takes a little thinking, but one can see that each symmetry is determined by its action on the faces. Let us consider symmetries of the tetrahedron that preserve orientation. For example, $(123)$ is such a symmetry; it rotates the tetrahedron relative to the plane on which the face 4 lies. On the other hand, $(12)$ is not such a symmetry. We conclude that the symmetries that preserve orientation are a subgroup of $S_4$ that is not equal to $S_4$. Let us call this subgroup $A_4$ (later on, we shall define the groups $A_n$ in general and our notation is consistent). Clearly $A_4$ contains \{1, (123), (132), (234), (243), (134), (143), (124), (142)\} and since it is closed under multiplication (being a subgroup) also $(132)(134) = (12)(34)$ and similarly, $(13)(24)$ and $(14)(23)$, are elements of $A_4$. We have already identified 12 elements in $A_4$ and since the order of $A_4$ divides the order of $S_4$, which is 24, $A_4$ is in fact equal to
\[\{1, (123), (132), (234), (243), (134), (143), (124), (142), (12)(34), (13)(24), (14)(23)\}.

Let us count how many colorings of the faces of the tetrahedron are there using 4 distinct colors, each once. The number of coloring is $4! = 24$ (choose for each face its color). No symmetry but the identity preserves a coloring and so by CFF we get that the number of colorings up to $A_4$ identifications is 2.

Suppose now we want to color with 2 colors, say red and blue, painting two faces red and two faces blue. The total number of colorings are \(\binom{4}{2} = 6\). In this case, a three cycle cannot fix a coloring, while each permutation of the type $(ab)(cd)$ fixes exactly two colorings (choose if the
faces \( a, b \) are both red or both blue. Therefore, the number of colorings up to symmetries is
\[
N = \frac{1}{12}(6 + 3 \times 2) = 1.
\]

28. Cauchy’s theorem: a wonderful proof

One application of group actions is to provide a simple proof of an important theorem in the theory of finite groups. Every other proof I know is very complicated.

**Theorem 28.1.** (Cauchy) Let \( G \) be a finite group of order \( n \) and let \( p \) be a prime dividing \( n \). Then \( G \) has an element of order \( p \).

**Proof.** \(^{14}\) Let \( S \) be the set consisting of \( p \)-tuples \((g_1, \ldots, g_p)\) of elements of \( G \), considered up to cyclic permutations. Thus if \( T \) is the set of \( p \)-tuples \((g_1, \ldots, g_p)\) of elements of \( G \), \( S \) is the set of orbits for the action of \( \mathbb{Z}/p\mathbb{Z} \) on \( T \) by cyclic shifts (\( \�� \)). One may therefore apply CFF and get
\[
|S| = \frac{n^p - n}{p} + n. \quad (\��)
\]

Note that \( n \nmid |S| \) (\( \�� \)).

Now define an action of \( G \) on \( S \). Given \( g \in G \) and \((g_1, \ldots, g_p) \in S \) we define
\[
g(g_1, \ldots, g_p) = (gg_1, \ldots, gg_p).
\]

This is a well defined action (\( \�� \)).

Since the order of \( G \) is \( n \), since \( n \nmid |S| \), and since \( S \) is a disjoint union of orbits of \( G \), there must be an orbit \( \text{Orb}(s) \) whose size is not \( n \). However, the size of an orbit is \( |G|/|\text{Stab}(s)| \), and we conclude that there must an element \((g_1, \ldots, g_p)\) in \( S \) with a non-trivial stabilizer. This means that for some \( g \in G \), such that \( g \neq e \), we have
\[
(gg_1, \ldots, gg_p) \text{ is equal to } (g_1, \ldots, g_p) \text{ up to a cyclic shift.}
\]

This means that for some \( i \) we have
\[
(gg_1, \ldots, gg_p) = (g_{i+1}, g_{i+2}, g_{i+3}, \ldots, g_p, g_1, g_2, \ldots, g_i).
\]

Therefore, \( gg_1 = g_{i+1}, g^2g_1 = gg_{i+1} = g_{2i+1}, \ldots, g^pg_1 = \cdots = g_{pi+1} = g_1 \) (we always read the indices mod \( p \)). That is, there exists \( g \neq e \) with
\[
g^p = e.
\]

Let \( k \) be the order of \( g \). Then \( 2 \leq k \leq p \). Write \( p = qk + r \) with \( 0 \leq r < k \). Then
\[
g^r = g^p \cdot (g^k)^{-q} = e \cdot (e)^q = e.
\]

\(^{14}\)Some details are left out. We use the symbol \( \�� \) to point out that a detail is left out and the interested person should fill it in.
Since \( k \) is the minimal positive power such that \( g^k = e \) we must have \( k \mid p \). This implies \( k = p \) and therefore we found an element of order \( p \) (namely, \( g \)).

\[ \square \]

29. The First Isomorphism Theorem for Groups

29.1. Normal subgroups.

Definition 29.1. Let \( G \) be a group and \( H \) a subgroup of \( G \). \( H \) is called a normal subgroup if for every \( g \in G \) we have

\[ gH = Hg. \]

Note that \( gH = Hg \) if and only if \( gHg^{-1} = H \), where \( gHg^{-1} = \{ghg^{-1} : h \in H\} \). Thus, we could also define a normal subgroup \( H \) to be a subgroup such that \( gHg^{-1} = H \) for all \( g \in G \), equivalently, \( \forall g \in G, \forall h \in H, ghg^{-1} \in H \).

Lemma 29.2. Let \( H \) be a subgroup of a group \( G \). Then \( H \) is normal if and only if

\[ gHg^{-1} \subseteq H, \quad \forall g \in G. \]

Proof. Clearly if \( H \) is normal, \( gHg^{-1} \subseteq H, \forall g \in G \). Suppose then that \( gHg^{-1} \subseteq H, \forall g \in G \).

Given \( g \in G \) we have then \( gHg^{-1} \subseteq H \) and also \( g^{-1}H(g^{-1})^{-1} \subseteq H \). The last inclusion is just \( g^{-1}Hg \subseteq H \), which is equivalent to \( H \subseteq gHg^{-1} \). We conclude that \( gHg^{-1} = H \). \( \square \)

Our main example of a normal subgroup is the kernel of a homomorphism.

Proposition 29.3. Let \( f : G \to H \) be a group homomorphism. Then \( \text{Ker}(f) \) is a normal subgroup of \( G \).

Proof. We proved already that \( \text{Ker}(f) \) is a subgroup. Let \( g \in G, h \in \text{Ker}(f) \); we need to show that \( ghg^{-1} \in \text{Ker}(f) \), that is \( f(ghg^{-1}) = e_H \). We calculate \( f(ghg^{-1}) = f(g)f(h)f(g^{-1}) = f(g)e_H f(g^{-1}) = f(g)f(g^{-1}) = f(g)f(g)^{-1} = e_H \). \( \square \)

Example 29.4. For any group \( G \), \( \{e_G\} \) and \( G \) are normal subgroups. If \( G \) is a commutative group, any subgroup of \( G \) is a normal subgroup.

29.2. Quotient groups. Similar to the construction of a quotient ring, we construct quotient groups.

Let \( G \) be a group and \( H \) a normal subgroup of \( G \). We let the quotient group \( G \) mod \( H \), denoted \( G/H \), be the collection of left cosets of \( H \). We define multiplication by

\[ (aH)(bH) = abH. \]

We claim that this is well defined, namely, if \( aH = a_1H, bH = b_1H \) then \( abH = a_1b_1H \). Indeed, we have \( a = a_1h \) for some \( h \in H \) and \( b = b_1h' \) for some \( h' \in H \). Also, \( hb_1 \in Hb_1 = b_1H \).
and so \( hb_1 = h''b_1 \) for some \( h'' \in H \). Then, \( abH = a_1hb_1h'H = a_1b_1h''h'H = a_1b_1H \) (if \( t \in H \) then \( tH = H \)).

We now verify the group axioms. We use the notation \( \bar{a} \) for \( aH \). Then the group law is

\[
\bar{a} \bar{b} = \overline{ab}.
\]

We have \((\bar{a} \bar{b})\bar{c} = \overline{ab} \overline{bc} = (ab)c = a\overline{bc} = \bar{a} \overline{bc} = \bar{a}(\bar{b} \bar{c})\). Thus, this is an associative operation.

We have \( \bar{a} \overline{e_G} = \overline{ae_G} = \bar{a} \) and \( \overline{e_G} \bar{a} = \overline{e_Ga} = \bar{a} \). So there is an identity element \( e_{G/H} \) and it is equal to \( \bar{e_G} = H \). Finally, \( \bar{a} a^{-1} = \overline{aa^{-1}} = \bar{e_G} = e_{G/H} \) and \( \overline{a^{-1}a} = \bar{a}^{-1}a = \bar{e_G} = e_{G/H} \). Thus, every \( \bar{a} \) is invertible and its inverse is \( \overline{a^{-1}} \) (that is, \( (aH)^{-1} = a^{-1}H \)).

### 29.3. The first isomorphism theorem.

**Theorem 29.5.** Let \( f : G \to H \) be a surjective group homomorphism. The canonical map \( \pi : G \to G/\text{Ker}(f) \) is a homomorphism with kernel \( \text{Ker}(f) \). There is an isomorphism \( F : G/\text{Ker}(f) \to H \), such that the following diagram commutes:

\[
\begin{array}{ccc}
G & \xrightarrow{f} & H \\
\pi \downarrow & & \downarrow F \\
G/\text{Ker}(f) & & 
\end{array}
\]

**Proof.** First we check that \( \pi : G \to G/\text{Ker}(f) \) is a homomorphism, where \( \pi(a) = \bar{a} = a\text{Ker}(f) \). Indeed, this is just the formula \( \overline{ab} = \bar{a} \bar{b} \). The kernel is \( \{a \in G : a\text{Ker}(f) = \text{Ker}(f)\} = \text{Ker}(f) \).

Let us define

\[
F : G/\text{Ker}(f) \to H, \quad F(\bar{a}) = f(a).
\]

This is well defined: if \( \bar{a} = \bar{b} \) then \( b^{-1}a \in \text{Ker}(f) \), so \( f(b) = f(b)f(b^{-1}a) = f(b(b^{-1}a)) = f(a) \).

Clearly \( F \circ \pi = f \).

\( F \) is a homomorphism: \( F(\bar{a} \bar{b}) = F(ab) = f(ab) = f(a)f(b) = F(\bar{a})F(\bar{b}) \). Furthermore, \( F \) is surjective, since given \( h \in H \) we may find \( a \in G \) such that \( f(a) = h \) and so \( F(\bar{a}) = h \). Finally, \( F \) is injective, because \( F(\bar{a}) = f(a) = e_H \) means that \( a \in \text{Ker}(f) \) so \( \bar{a} = e_{G/H} \). \( \square \)

**Example 29.6.** Let \( \mathbb{F} \) be a field. Recall the group of matrices \( \text{GL}_2(\mathbb{F}) \),

\[
\text{GL}_2(\mathbb{F}) = \left\{ M = \begin{pmatrix} a & b \\ c & d \end{pmatrix} : a, b, c, d \in \mathbb{F}, \det(M) = ad - bc \neq 0 \right\}.
\]

We have also noted that the determinant is multiplicative

\[
\det(MN) = \det(M) \det(N).
\]

We may now view this fact as saying that the function

\[
\det : \text{GL}_2(\mathbb{F}) \to \mathbb{F}^*,
\]
is a group homomorphism. It is a surjective group homomorphism, because given any \( a \in \mathbb{F}^\times \) the matrix \[
abla = \begin{pmatrix} a & 0 \\ 0 & 1 \end{pmatrix}
\] has determinant \( a \). The kernel is called \( \text{SL}_2(\mathbb{F}) \), it is equal to the matrices with determinant 1. It is a normal subgroup of \( \text{GL}_2(\mathbb{F}) \) and by the first isomorphism theorem \( \text{GL}_2(\mathbb{F})/\text{SL}_2(\mathbb{F}) \cong \mathbb{F}^\times \).

**Example 29.7.** The homomorphic images of \( S_3 \). We wish to identify all the homomorphic images of \( S_3 \). If \( f : S_3 \to G \) is a group homomorphism then \( \text{Ker}(f) \) is a normal subgroup of \( S_3 \). We begin therefore by finding all normal subgroups of \( S_3 \).

We know that every nontrivial subgroup of \( S_3 \) is of the form \( \langle (ij) \rangle \) for some transposition \( (ij) \) or the subgroup \( A_3 := \langle (123) \rangle \). That there are no other subgroups follows from the following observation: if \( H \subset K \subset G \) are groups and \( G \) is finite, then \( |G|/|K| \) divides \( |G|/|H| \), because the quotient is \( |K|/|H| \). In our situation, for a non-trivial subgroup \( H \) we have \( |S_3|/|H| \) is either 2 or 3 and those are prime. It follows that either \( |K| = |G| \) or \( |K|/|H| \) and so that either \( K = G \) or \( K = H \).

The subgroups of order 2 are not normal. For example, \( (13)(12)(13)^{-1} = (13)(12)(13) = (23) \), which shows that \( \{1, (12)\} \) is not normal, etc. On the other hand, the subgroup \( A_3 := \{1, (123), (132)\} \) is normal. This follows from it being of index 2 (see assignments); another argument appears below. Since \( S_3/A_3 \) has order 2, it must be isomorphic to \( \mathbb{Z}/2\mathbb{Z} \).

We conclude that there are three options:

1. \( \text{Ker}(f) = \{1\} \). In this case, \( S_3 \) is isomorphic to its image.
2. \( \text{Ker}(f) = S_3 \). In this case \( S_3/\text{Ker}(f) = S_3/S_3 \cong \{1\} \) is the trivial group.
3. \( \text{Ker}(f) = A_3 \). In this case \( S_3/A_3 \) is a group of 2 elements, obviously cyclic. Thus \( S_3/A_3 \cong \mathbb{Z}/2\mathbb{Z} \).

29.4. **Groups of low order.**

29.4.1. **Groups of order 1.** There is a unique group of order 1, up to isomorphism. It consists of its identity element alone. There is only one way to define a homomorphism between two groups of order 1 and it is an isomorphism.

29.4.2. **Groups of order 2, 3, 5, 7.** Recall that we proved that every group \( G \) of prime order is cyclic, and, in fact, any non-trivial element is a generator. This implies that any subgroup of \( G \) different from \( \{e_G\} \) is equal to \( G \). We also proved that any two cyclic groups having the same order are isomorphic. We therefore conclude:

**Corollary 29.8.** Every group \( G \) of prime order \( p \) is isomorphic to \( \mathbb{Z}/p\mathbb{Z} \); it has no subgroups apart from the trivial subgroups \( \{e_G\}, G \).

In particular, this corollary applies to groups of order 2, 3, 5, 7.
29.4.3. Groups of order 4. Let $G$ be a group of order 4.

**First case:** $G$ is cyclic.

In this case we have $G \cong \mathbb{Z}/4\mathbb{Z}$. Its subgroups are \{0\}, $\mathbb{Z}/4\mathbb{Z}$ and $H = \langle 2 \rangle = \{0, 2\}$. There are no other subgroups because if a subgroup $J$ contains an element $g$ it contains the cyclic subgroup generated by $g$. In our case, the elements 1 and 3 are generators, so any subgroup not equal to $G$ is contained in \{0, 2\}.

Since $G$ is abelian, $H$ is normal. $G/H$ has order $|G|/|H| = 4/2 = 2$ and so $G/H \cong \mathbb{Z}/2\mathbb{Z}$.

**Second case.** $G$ is not cyclic.

**Claim:** Every element of $G$ different from $e_G$ has order 2.

**Proof:** we have $\text{ord}(g) = |\langle g \rangle|$ and it divides $|G|$. So, in our case, $\text{ord}(g) = 1, 2$ or 4. If $\text{ord}(g) = 4$, we get that $G$ is cyclic and if $\text{ord}(g) = 1$ then $g = e_G$. Thus, we must have $\text{ord}(g) = 2$.

**Claim:** Let $G$ be a group in which every element different from the identity has order 2. Then $G$ is commutative.

**Proof:** Note first that if $a \in G$ has order 2 (or is the identity) then $aa = e_G$ and so $a^{-1} = a$. Now, we need to show that for every $a, b \in G$ we have $ab = ba$. But this is equivalent to $ab = b^{-1}a^{-1}$.

Multiply both sides by $ab$ and we see that we need to prove that $abab = e_G$. But, $abab = (ab)^2$ and so is equal to $e_G$, by assumption.

One example of a group of order 4 satisfying all these properties is $\mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2\mathbb{Z}$. We claim that $G \cong \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2\mathbb{Z}$. Pick two distinct elements $g_1, g_2$ of $G$ that aren’t the identity either. Define a map

$$f : \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2\mathbb{Z} \to G, \quad f(a, b) = g_1^a g_2^b.$$ 

This is well defined: if $(a, b) = (a', b')$ then $a = a' + 2c, b = b' + 2d$ and we get $f(a, b) = g_1^a g_2^b = g_1^{a'}(g_1^2)^c (g_2^2)^d g_1^{b'} g_2^b = f(a', b')$. The map is also a homomorphism: $f((a_1, b_1) + (a_2, b_2)) = f(a_1 + a_2, b_1 + b_2) = g_1^{a_1 + a_2} g_2^{b_1 + b_2} = g_1^{a_1} g_1^{a_2} g_2^{b_1} g_2^{b_2}$. Because $G$ is commutative we can rewrite this as $f(a_1 + a_2, b_1 + b_2) = g_1^{a_1} g_2^{b_1} g_1^{a_2} g_2^{b_2} = f(a_1, b_1) \cdot f(a_2, b_2)$.

The image of $f$ is a subgroup with at least 3 elements, namely, $e_G, g_1, g_2$. By Lagrange the image then must be $G$. It follows that $f$ is surjective and so is also injective.

The non-trivial subgroups of $\mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/2\mathbb{Z}$ are all cyclic. They are \{(0, 0), (0, 1)\}, \{(0, 0), (1, 0)\} and \{(0, 0), (1, 1)\}. Since the group is commutative they are all normal and the quotient in every case has order 2, hence isomorphic to $\mathbb{Z}/2\mathbb{Z}$.

29.4.4. Groups of order 6. We know three candidates already $\mathbb{Z}/6\mathbb{Z}, \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/3\mathbb{Z}$ and $S_3$. Now, in fact, $\mathbb{Z}/6\mathbb{Z} \cong \mathbb{Z}/2\mathbb{Z} \times \mathbb{Z}/3\mathbb{Z}$ (say by CRT). And since $S_3$ is not commutative it is not isomorphic to $\mathbb{Z}/6\mathbb{Z}$. In fact, every group of order 6 is isomorphic to either $\mathbb{Z}/6\mathbb{Z}$ or $S_3$. We don’t prove it here.
The subgroups of $\mathbb{Z}/6\mathbb{Z}$: Let $n$ be a positive integer. We have a surjective group homomorphism $\pi : \mathbb{Z} \to \mathbb{Z}/n\mathbb{Z}$. Similar to the situation with rings one can show that this gives a bijection between subgroups $H$ of $\mathbb{Z}$ that contain $n\mathbb{Z}$ and subgroups $K$ of $\mathbb{Z}/n\mathbb{Z}$. The bijection is given by

$$H \mapsto \pi(H), \quad K \mapsto \pi^{-1}(K).$$

The subgroups of $\mathbb{Z}$ are all cyclic, having the form $n\mathbb{Z}$ for some $n$ (same proof as for ideals, really). We thus conclude that the subgroups of $\mathbb{Z}/n\mathbb{Z}$ are cyclic and generated by the elements $m$ such that $m|n$. Thus, for $n = 6$ we find the cyclic subgroups generated by $1, 2, 3, 6$. Those are $\mathbb{Z}/6\mathbb{Z}, \{0, 2, 4\}, \{0, 3\}, \{0\}$. They are all normal and the quotients are isomorphic respectively to $\{0\}, \mathbb{Z}/2\mathbb{Z}, \mathbb{Z}/3\mathbb{Z}, \mathbb{Z}/6\mathbb{Z}$.

The subgroups of $S_3$: Those were classified above.

29.5. **Odds and evens.** Let $n \geq 2$ be an integer. One can show that there is a way to assign a sign, $\pm 1$, to any permutation in $S_n$ such that the following properties hold:

- $\text{sgn}(\sigma\tau) = \text{sgn}(\sigma) \cdot \text{sgn}(\tau)$.
- $\text{sgn}((ij)) = -1$ for $i \neq j$.

We do not prove that here, but we shall prove that next term in MATH 251. Note that since any permutation is a product of transpositions, the two properties together determine the sign of any permutation. Here are some examples: $\text{sgn}((12)) = -1, \text{sgn}((123)) = \text{sgn}((13)(12)) = \text{sgn}(13) \cdot \text{sgn}(12) = 1, \text{sgn}((1234)) = \text{sgn}((14)(13)(12)) = -1^3 = -1$.

The property $\text{sgn}(\sigma\tau) = \text{sgn}(\sigma) \cdot \text{sgn}(\tau)$ could be phrased as saying that the function $\text{sgn} : S_n \to \{\pm 1\}$ is a surjective group homomorphism. We define $A_n$ as the kernel of the homomorphism $\text{sgn}$. It is called the alternating group on $n$ letters and its elements are called even permutations. The elements of $S_n \setminus A_n$ are called odd permutations. The group $A_n$ is a normal subgroup of $S_n$, being a kernel of a homomorphism. Its cardinality is $n!/2$. Here are some examples:

- $A_2 = \{1\}$;
- $A_3 = \{1, (123), (132)\}$;
- $A_4 = \{1, (12)(34), (13)(24), (14)(23), (123), (132), (234), (243), (124), (142), (134), (143)\}$.

(Easy to check those are distinct 12 even permutations, so the list must be equal to $A_4$).

29.6. **Odds and Ends.**

**Example 29.9.** We prove that in $\mathbb{Z}/p\mathbb{Z}$ any element is a sum of two squares.

Clearly this holds for $p = 2$, so we assume $p > 2$. To begin with, $\mathbb{Z}/p\mathbb{Z}^\times = \{1, \ldots, p-1\}$ is a group under multiplication; it has $p - 1$ elements. Consider the homomorphism

$$\text{sq} : \mathbb{Z}/p\mathbb{Z}^\times \to \mathbb{Z}/p\mathbb{Z}^\times, \quad \text{sq}(x) = x^2.$$
Let $H$ be its image – a subgroup of $\mathbb{Z}/p\mathbb{Z}^\times$. The kernel of $sq$ is the solutions to $x^2 = 1$, which are precisely $\pm 1$. Note that $1 \neq -1$. It follows that $H \cong \mathbb{Z}/p\mathbb{Z}^\times/\{\pm 1\}$ is a group with $(p-1)/2$ elements consisting precisely of the non-zero congruence classes that are squares.

Let $H^* = H \cup \{0\}$; it is a subset of $\mathbb{Z}/p\mathbb{Z}$ with $(p+1)/2$ elements consisting of all squares.

Let $a \in \mathbb{Z}/p\mathbb{Z}$ then the two sets $H^*$ and $a - H^* := \{a - h : h \in H^*\}$ have size $(p+1)/2$ and so must intersect (because $\mathbb{Z}/p\mathbb{Z}$ has $p < 2 \cdot \frac{p+1}{2}$ elements). That is, there are two squares $x^2, y^2$ such that $a - x^2 = y^2$ and so $a = x^2 + y^2$.

We next tie together the notions of homomorphism and group action.

**Lemma 29.10.** Let $G$ be a group and $T$ a non-empty set. To give an action of $G$ on $T$ is equivalent to giving a homomorphism $\rho : G \to S_T$.

**Proof.** Suppose that we are given an action of $G$ on $S$. Pick an element $g \in G$. We claim that the function

$$T \to T, \quad t \mapsto g \ast t,$$

is a permutation of $T$. Indeed, if $gt_1 = gt_2$ then $g^{-1}(gt_1) = g^{-1}(gt_2)$, so $(g^{-1}g)t_1 = (g^{-1}g)t_2$; that is, $et_1 = et_2$ and so $t_1 = t_2$. Also, given $t \in T$ we have $g(g^{-1}t) = (gg^{-1})t = et = t$, showing surjectivity. Let us denote then this function by $\rho(g)$, $\rho(g)t = gt$. We have a function

$$\rho : G \to S_T.$$

We claim that this function is a homomorphism. We need to show that $\rho(g_1g_2) = \rho(g_1) \circ \rho(g_1)$, i.e., that for every $t \in T$ we have $\rho(g_1g_2)(t) = (\rho(g_1) \circ \rho(g_2))(t)$. Indeed, $\rho(g_1g_2)t = (g_1g_2)t = g_1(g_2)t = g_1(\rho(g_2)t) = \rho(g_1)(\rho(g_2)(t)) = (\rho(g_1) \circ \rho(g_2))(t)$.

Conversely, suppose that

$$\rho : G \to S_T$$

is a group homomorphism. Define an action of $G$ on $S$ by

$$g \ast t := \rho(g)(t).$$

We claim this is a group action. Since $\rho$ is a homomorphism we have $\rho(e) = \text{Id}_T$ and so $e \ast t = \rho(e)(t) = \text{Id}_T(t) = t$. Now, $g_1 \ast (g_2 \ast t) = \rho(g_1)(\rho(g_2)(t)) = (\rho(g_1) \circ \rho(g_2))(t) = \rho(g_1g_2)(t) = (g_1g_2) \ast t$. \hfill $\square$
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