# Not Always Buried Deep 

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Dedicated to the memory of Arnold Ephraim Ross (1906-2002).

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## Foreword

The gold in 'them there hills' is not always buried deep. Much of it is within easy reach. Some of it is right on the surface to be picked up by any searcher with a keen eye for detail and an eagerness to explore. As in any treasure hunt, the involvement grows as the hunt proceeds and each success whether small or great adds the fuel of excitement to the exploration. - A. E. Ross

Number theory is one of the few areas of mathematics where problems of substantial interest can be described to someone possessing scant mathematical background. It sometimes proves to be the case that a problem which is simple to state requires for its resolution considerable mathematical preparation; e.g., this appears to be the case for Fermat's conjecture regarding integer solutions to the equation $x^{n}+y^{n}=z^{n}$. But this is by no means a universal phenomenon; many engaging problems can be successfully attacked with little more than one's "mathematical bare hands". In this case one says that the problem can be solved in an elementary way (even though the elementary solution may be far from simple). Such elementary methods and the problems to which they apply are the subject of this book.

Because of the nature of the material, very little is required in terms of prerequisites: The reader is expected to have prior familiarity with number theory at the level of an undergraduate course. The nececssary background can be gleaned from any number of excellent texts, such as Sierpiński's charmingly discursive Elementary Theory of Numbers or LeVeque's lucid and methodical Fundamentals of Number Theory. Apart from this, a rigorous course in calculus, some facility with manipulation of estimates (in
particular, big-Oh and little-oh notation), and a first course in modern algebra (covering groups, rings, and fields) should suffice for the majority of the text. A course in complex variables is not required, provided that the reader is willing to overlook some motivational remarks made in Chapter 7.

Rather than attempt a comprehensive account of elementary methods in number theory, I have focused on topics that I find particularly attractive and accessible:

- Chapters $1,3,4$, and 7 collectively provide an overview of prime number theory, starting from the infinitude of the primes, moving through the elementary estimates of Chebyshev and Mertens, then the theorem of Dirichlet on primes in prescribed arithmetic progressions, and culminating in an elementary proof of the prime number theorem.
- Chapter 2 contains a discussion of Gauss's arithmetic theory of the roots of unity (cyclotomy), which was first presented in the final section of his Disquisitiones Arithmeticae. After developing this theory to the extent required to prove Gauss's characterization of constructible regular polygons, we give a cyclotomic proof of the quadratic reciprocity law, followed by a detailed account of a little-known cubic reciprocity law due to Jacobi.
- Chapter 5 is a 12 -page interlude containing Dress's proof of the following result conjectured by Waring in 1770 and established by Hilbert in 1909: For each fixed integer $k \geq 2$, every natural number can be expressed as the sum of a bounded number of nonnegative $k$ th powers, where the bound depends only on $k$.
- Chapter 6 is an introduction to combinatorial sieve methods, which were introduced by Brun in the early twentieth century. The bestknown consequence of Brun's method is that if one sums the reciprocals of each prime appearing in a twin prime pair $p, p+2$, then the answer is finite. Our treatment of sieve methods is robust enough to establish not only this and other comparable 'upper bound' results, but also Brun's deeper "lower bound" results. For example, we prove that there are infinitely many $n$ for which both $n$ and $n+2$ have at most 7 prime factors, counted with multiplicity.
- Chapter 8 summarizes what is known at present about perfect numbers, numbers which are the sum of their proper divisors.

At the end of each chapter (excepting the interlude) I have included several nonroutine exercises. Many are based on articles from the mathematical literature, including both research journals and expository publications like the American Mathematical Monthly. Here, as throughout the text, I have
made a conscious effort to document original sources and thus encourage conformance to Abel's advice to "read the masters".

While the study of elementary methods in number theory is one of the most accessible branches of mathematics, the lack of suitable textbooks has been a repellent to potential students. It is hoped that this modest contribution will help to reverse this injustice.

Paul Pollack

## Notation

While most of our notation is standard and should be familiar from an introductory course in number theory, a few of our conventions deserve explicit mention: The set $\mathbf{N}$ of natural numbers is the set $\{1,2,3,4, \ldots\}$. Thus 0 is not considered a natural number. Also, if $n \in \mathbf{N}$, we write " $\tau(n)$ " (instead of " $d(n)$ ") for the number of divisors of $n$. This is simply to avoid awkward expressions like " $d(d)$ " for the number of divisors of the natural number $d$. Throughout the book, we reserve the letter $p$ for a prime variable.

We remind the reader that " $A=O(B)$ " indicates that $|A| \leq c|B|$ for some constant $c>0$ (called the implied constant); an equivalent notation is " $A \ll B$ ". The notation " $A \gg B$ " means $B \ll A$, and we write " $A \asymp B$ " if both $A \ll B$ and $A \gg B$. If $A$ and $B$ are functions of a single real variable $x$, we often speak of an estimate of this kind holding as " $x \rightarrow a$ " (where $a$ belongs to the two-point compactification $\mathbf{R} \cup\{ \pm \infty\}$ of $\mathbf{R}$ ) to mean that the estimate is valid on some deleted neighborhood of $a$. Subscripts on any of these symbols indicate parameters on which the implied constants (and, if applicable, the deleted neighborhoods) may depend. The notation " $A \sim B$ " means $A / B \rightarrow 1$ while " $A=o(B)$ " means $A / B \rightarrow 0$; here subscripts indicate parameters on which the rate of convergence may depend.

If $S$ is a subset of the natural numbers $\mathbf{N}$, the (asymptotic, or natural) density of $S$ is defined as the limit

$$
\lim _{x \rightarrow \infty} \frac{1}{x} \#\{n \in S: n \leq x\},
$$

provided that this limit exists. The lower density and upper density of $S$ are defined similarly, with liminf and lim sup replacing lim (respectively). We say that a statement holds for almost all natural numbers $n$ if it holds on a subset of $\mathbf{N}$ of density 1 .

If $f$ and $G$ are defined on a closed interval $[a, b] \subset \mathbf{R}$, with $f^{\prime}$ piecewise continuous there, we define

$$
\begin{equation*}
\int_{a}^{b} f(t) d G(t):=G(b) f(b)-G(a) f(a)-\int_{a}^{b} f^{\prime}(t) G(t) d t \tag{0.1}
\end{equation*}
$$

provided that the right-hand integral exists. (Experts will recognize the right-hand side as the formula for integration by parts for the RiemannStieltjes integral, but defining the left-hand side in this manner allows us to avoid assuming any knowledge of Riemann-Stieltjes integration.) We will often apply partial summation in the following form, which is straightforward to verify directly: Suppose that $a$ and $b$ are real numbers with $a \leq b$ and that we are given complex numbers $a_{n}$ for all natural numbers $n$ with $a<n \leq b$. Put $S(t):=\sum_{a<n \leq t} a_{n}$. If $f^{\prime}$ is piecewise continuous on $[a, b]$, then

$$
\sum_{a<n \leq b} a_{n} f(n)=\int_{a}^{b} f(t) d S(t) .
$$

In order to paint an accurate portrait of the mathematical landscape without straying off point, it has been necessary on occasion to state certain theorems without proof; such results are marked with a star ( $\boldsymbol{\star})$. For some of these results, proofs are sketched in the corresponding chapter exercises.

## Acknowledgements

There are many people without whom this book could not have been written and many others without whom this book would not be worth reading.

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My family - my father Lawrence, my mother Lolita, and my brother Michael - has done so much for me over the years that it would be impossible (and inappropriate!) for me to express the extent of my appreciation in this brief space. Another friend for whom I am grateful beyond words is Susan Roth, who for the last decade has accompanied me on many of my (mis)adventures in genre television.

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Finally, I would like to thank the American Mathematical Society, especially Ed Dunne, Cristin Zanella, and Luann Cole, for their encouragement of this project at every stage.

## Elementary Prime Number Theory, I

Prime numbers are more than any assigned multitude of prime numbers. - Euclid

No prime minister is a prime number - A. Plantinga

## 1. Introduction

Recall that a natural number larger than 1 is called prime if its only positive divisors are 1 and itself, and composite otherwise. The sequence of primes begins

$$
2,3,5,7,11,13,17,19,23,29,31,37,41,43, \ldots
$$

Few topics in number theory attract more attention, popular or professional, than the theory of prime numbers. It is not hard to see why. The study of the distribution of the primes possesses in abundance the very features that draw so many of us to mathematics in the first place: intrinsic beauty, accessible points of entry, and a lingering sense of mystery embodied in numerous unpretentious but infuriatingly obstinate open problems.

Put

$$
\pi(x):=\#\{p \leq x: p \text { prime }\} .
$$

Prime number theory begins with the following famous theorem from antiquity:

Theorem 1.1. There are infinitely many primes, i.e., $\pi(x) \rightarrow \infty$ as $x \rightarrow \infty$.

The first half of this chapter is a survey of the many proofs that have been given for Theorem 1.1. The second half of this chapter is devoted to the theme of prime-producing formulas and the occurrence of primes in various natural sequences.

## 2. Euclid and his imitators

We begin with the classic proof from Euclid's Elements (circa 300 BC ):
Proof. Suppose that $p_{1}, p_{2}, \ldots, p_{k}$ is any finite list of primes. Let $P$ denote the product of the $p_{i}$ and consider the integer $P+1$. Since $P+1 \equiv 1\left(\bmod p_{i}\right)$ for each $1 \leq i \leq k$, none of the $p_{i}$ divide $P+1$. But since $P+1>1$, it must have some prime divisor $p$. It follows that there is always a prime missing from any finite list, or, as Euclid put it, "prime numbers are more than any assigned multitude of primes."

There are many trivial variants; for instance, we can easily show that for every integer $m$ there is a prime $p>m$ by taking $p$ to be any prime divisor of $m!+1$.

In this section we collect several Euclidean proofs for Theorem 1.1. All of these start with a finite list of primes and then produce an integer $>1$ that is coprime to every prime on the list. Stieltjes's proof is typical:

Stieltjes's proof, 1890. Suppose that $p_{1}, \ldots, p_{k}$ is a finite list of distinct primes with product $P$ and let $P=A B$ be any decomposition of $P$ into two positive factors. Suppose that $p$ is one of the $p_{i}$. Then $p \mid A B$, so that either $p \mid A$ or $p \mid B$. If $p$ divides both $A$ and $B$, then $p^{2}$ divides $P$, which is false. Consequently, $p$ divides exactly one of $A$ and $B$. It follows that $p \nmid A+B$. So $A+B$ is divisible by none of the $p_{i}$; but as $A+B \geq 2$, it has some prime divisor. So again we have discovered a prime not on our original list.

Euler's second proof (published posthumously). This proof is based on the multiplicativity of the Euler totient function: Let $p_{1}, \ldots, p_{k}$ be a list of distinct primes with product $P$. By said multiplicativity,

$$
\varphi(P)=\prod_{i=1}^{k}\left(p_{i}-1\right) \geq 2^{k-1} \geq 2
$$

provided that our list contains at least two primes (as we may assume). It follows that there is an integer in the interval $[2, P]$ that is coprime to $P$; but such an integer has a prime factor distinct from all of the $p_{i}$.

Proof of Braun (1897), Métrod (1917). Let $p_{1}, \ldots, p_{k}$ be a list of $k \geq$ 2 distinct primes and let $P=p_{1} p_{2} \cdots p_{k}$. Consider the integer

$$
N:=P / p_{1}+P / p_{2}+\cdots+P / p_{k} .
$$

For each $1 \leq i \leq k$, we have

$$
N \equiv P / p_{i}=\prod_{j \neq i} p_{j} \not \equiv 0\left(\bmod p_{i}\right)
$$

so that $N$ is divisible by none of the $p_{i}$. But $N \geq 2$, and so it must possess a prime factor not on our list.

## 3. Coprime integer sequences

Suppose we know an infinite sequence of pairwise relatively prime positive integers

$$
2 \leq n_{1}<n_{2}<\cdots .
$$

Then we may define a sequence of primes $p_{i}$ by selecting arbitrarily a prime divisor of the corresponding $n_{i}$; the terms of this sequence are pairwise distinct because the $n_{i}$ are pairwise coprime.

If we can exhibit such a sequence of $n_{i}$ without invoking the infinitude of the primes, then we have a further proof of Theorem 1.1. An argument of this nature was given by Goldbach:

Proof (Goldbach). Let $n_{1}=3$, and for $i>1$ inductively define

$$
n_{i}=2+\prod_{1 \leq j<i} n_{j} .
$$

The following assertions are all easily verified in succession:
(i) Each $n_{i}$ is odd.
(ii) When $j>i$, we have $n_{j} \equiv 2 \bmod n_{i}$.
(iii) We have $\operatorname{gcd}\left(n_{i}, n_{j}\right)=1$ for $i \neq j$.

Theorem 1.1 now follows from the above remarks.
A straightforward induction shows that

$$
\begin{equation*}
n_{i}=2^{2^{i-1}}+1, \tag{1.1}
\end{equation*}
$$

and this is how Goldbach presented the proof.
Before proceeding, we pause to note that the above proof implies more than simply the infinitude of the primes. First, it gives us an upper bound for the $n$th prime, $2^{2^{n-1}}+1$; this translates into a lower bound of the shape

$$
\pi(x) \gg \log \log x \quad(x \rightarrow \infty)
$$

Second, it may be used to prove that certain arithmetic progressions contain infinitely many primes. To see this, suppose that $p \mid n_{i}$ and note that by (1.1), we have

$$
2^{2^{i-1}} \equiv-1(\bmod p), \quad \text { so that } 2^{2^{i}} \equiv\left(2^{2^{i-1}}\right)^{2} \equiv 1(\bmod p) .
$$

Hence the order of 2 modulo $p$ is precisely $2^{i}$. Thus $2^{i} \mid(\mathbf{Z} / p \mathbf{Z})^{\times}=p-1$, so that $p \equiv 1\left(\bmod 2^{i}\right)$. As a consequence, for any fixed $k$, there are infinitely many primes $p \equiv 1\left(\bmod 2^{k}\right)$ : choose a prime $p_{i}$ dividing $n_{i}$ for each $i \geq k$. In $\S 9.1$ we will prove the more general result that for each $m \geq 1$, there are infinitely many primes $p \equiv 1(\bmod m)$.

A related method of proving the infinitude of the primes is as follows: Let $a_{1}<a_{2}<a_{3}<\cdots$ be a sequence of positive integers with the property that

$$
\operatorname{gcd}(i, j)=1 \quad \Longrightarrow \quad \operatorname{gcd}\left(a_{i}, a_{j}\right)=1
$$

Moreover, suppose that for some prime $p$, the integer $a_{p}$ has at least two distinct prime divisors. Then if $p_{1}, \ldots, p_{k}$ were a list of all the primes, the integer

$$
a_{p_{1}} a_{p_{2}} \cdots a_{p_{k}}
$$

would possess at least $k+1$ prime factors: indeed, each factor exceeds 1 , the factors are pairwise relatively prime, and one of the factors is divisible by two distinct primes. So there are $k+1>k$ primes, a contradiction.

It remains to construct such a sequence. We leave to the reader the easy exercise of showing that $a_{n}=2^{n}-1$ has the desired properties (note that $a_{11}=23 \cdot 89$ ). The original version of this argument, where $a_{n}$ is instead chosen as the $n$th Fibonacci number, is due to Wunderlich [Wun65]. The generalization presented here is that of Hemminger [Hem66].

Saidak [Sai06] has recently given a very simple argument making use of coprimality. Start with a natural number $n>1$. Because $n$ and $n+1$ are coprime, the number $N_{2}:=n(n+1)$ must have at least two distinct prime factors. By the same reasoning,

$$
N_{3}:=N_{2}\left(N_{2}+1\right)=n(n+1)(n(n+1)+1)
$$

must have at least three distinct prime factors. In general, having constructed $N_{j}$ with at least $j$ different prime factors, the number $N_{j+1}:=$ $N_{j}\left(N_{j}+1\right)$ must have at least $j+1$.

## 4. The Euler-Riemann zeta function

For complex numbers $s$ with real part greater than 1, define the zeta function by putting

$$
\zeta(s):=\sum_{n=1}^{\infty} \frac{1}{n^{s}} .
$$

(The condition that $\Re(s)>1$ guarantees convergence of the series.) In the analytic approach to prime number theory, this function occupies a central position. Because of this text's emphasis on elementary methods, the zeta function will not play a large role for us, but it should be stressed that in many of the deeper investigations into the distribution of primes, the zeta function is an indispensable tool.

Riemann introduced the study of $\zeta(s)$ as a function of a complex variable in an 1859 memoir on the distribution of primes [Rie59]. But the connection between the zeta function and prime number theory goes back earlier. Over a hundred years prior to Riemann's study, Euler had looked at the same series for real $s$ and had shown that [Eul37, Theorema 8]

$$
\begin{equation*}
\sum_{n=1}^{\infty} \frac{1}{n^{s}}=\prod_{p} \frac{1}{1-\frac{1}{p^{s}}} \quad(s>1) . \tag{1.2}
\end{equation*}
$$

This is often called an analytic statement of unique factorization. To see why, notice that formally (i.e., disregarding matters of convergence)

$$
\prod_{p}\left(1+\frac{1}{p^{s}}+\frac{1}{p^{2 s}}+\cdots\right)=\sum_{n=1}^{\infty} \frac{a_{n}}{n^{s}},
$$

where $a_{n}$ counts the number of factorizations of $n$ into prime powers. Thus unique factorization, the statement that $a_{n}=1$ for all $n$, is equivalent to the statement that (1.2) holds as a formal product of Dirichlet series. ${ }^{1}$ This, in turn, is equivalent to the validity of (1.2) for all real $s>1$ (or even a sequence of $s$ tending to $\infty$ ) by a standard result in the theory of Dirichlet series (see, e.g., [Apo76, Theorem 11.3]).

Euler's product expansion of the zeta function is the first example of what is now called an Euler factorization. We now prove (following [Hua82]) a theorem giving general conditions for the validity of such factorizations.

Theorem 1.2 (Euler factorizations). Let $f$ be a multiplicative function. Then

$$
\begin{equation*}
\sum_{n=1}^{\infty} f(n)=\prod_{p}\left(1+f(p)+f\left(p^{2}\right)+\cdots\right) \tag{1.3}
\end{equation*}
$$

if either of the following two conditions holds:
(i) $\sum_{n=1}^{\infty}|f(n)|$ converges.
(ii) $\prod_{p}\left(1+|f(p)|+\left|f\left(p^{2}\right)\right|+\cdots\right)$ converges.

[^0]Remark. Without imposing a condition such as (i) or (ii), it is possible for either the series or the product in (1.3) to converge while the other diverges, or for both to converge without being equal. See [Win43, §15] for explicit examples.

If $f$ is not merely multiplicative but completely multiplicative, then the factors in (1.3) form a geometric series whose convergence is implied by either of the above conditions. Thus we have the following consequence:

Corollary 1.3. Let $f$ be a completely multiplicative function. Then

$$
\sum_{n=1}^{\infty} f(n)=\prod_{p} \frac{1}{1-f(p)}
$$

subject to either of the two convergence criteria of Theorem 1.2.
The factorization (1.2) of the zeta function is immediate from this corollary: One takes $f(n)=1 / n^{s}$ and observes that for $s>1$, condition (i) holds (for example) by the integral test.

Proof of Theorem 1.2. Suppose that condition (i) holds and set $S_{0}:=$ $\sum_{n=1}^{\infty}|f(n)|$. For each prime $p$, the series $\sum_{k=0}^{\infty} f\left(p^{k}\right)$ converges absolutely, since $\sum_{k=0}^{\infty}\left|f\left(p^{k}\right)\right| \leq S_{0}$. Therefore

$$
P(x)=\prod_{p \leq x}\left(1+f(p)+f\left(p^{2}\right)+\cdots\right)
$$

is a finite product of absolutely convergent series. It follows that

$$
P(x)=\sum_{n: p \mid n=p \leq x} f(n)
$$

If we now set $S=\sum_{n=1}^{\infty} f(n)$ (which converges absolutely), we have

$$
S-P(x)=\sum_{n: p \mid n \text { for some } p>x} f(n),
$$

which shows

$$
|S-P(x)| \leq \sum_{n>x}|f(n)| \rightarrow 0
$$

as $x \rightarrow \infty$. Thus $P(x) \rightarrow S$ as $x \rightarrow \infty$, which is the assertion of (1.3).
Now suppose that (ii) holds. We shall show that (i) holds as well, so that the theorem follows from what we have just done. To see this, let

$$
P_{0}=\prod_{p}\left(1+|f(p)|+\left|f\left(p^{2}\right)\right|+\cdots\right),
$$

and let

$$
\begin{aligned}
P_{0}(x) & :=\prod_{p \leq x}\left(1+|f(p)|+\left|f\left(p^{2}\right)\right|+\cdots\right) \\
& =\sum_{n: p \mid n \Rightarrow p \leq x}|f(n)| \geq \sum_{n \leq x}|f(n)| .
\end{aligned}
$$

Since $P_{0}(x) \leq P_{0}$ for all $x$, the partial sums $\sum_{n \leq x}|f(n)|$ form a bounded increasing sequence. Thus $\sum|f(n)|$ converges, proving (i).

We can now present Euler's first proof of the infinitude of the primes.
Euler's first proof of Theorem 1.1. Let $f$ be defined by $f(n)=1 / n$ for every $n$. Assuming that there are only finitely many primes, condition (ii) of Theorem 1.3 is trivially satisfied, as the product in question only has finitely many terms. It follows that

$$
\sum_{n=1}^{\infty} \frac{1}{n}=\prod_{p}\left(1+\frac{1}{p}+\frac{1}{p^{2}}+\cdots\right)<\infty
$$

in contradiction with the well-known divergence of the harmonic series.
As pointed out by Euler, this proof gives a much stronger result than that asserted in Theorem 1.1.

Theorem 1.4. The series $\sum \frac{1}{p}$ diverges, where the sum extends over all primes $p$.

Proof. Suppose not and let $C=\sum 1 / p$. As in the last proof, we take $f(n)=1 / n$ and apply Theorem 1.2. Let us check that condition (ii) of that theorem holds here. First, notice that

$$
\prod_{p \leq x}\left(1+\frac{1}{p}+\frac{1}{p^{2}}+\cdots\right)=\prod_{p \leq x} \frac{1}{1-\frac{1}{p}}=\prod_{p \leq x}\left(1+\frac{1}{p-1}\right) \leq \prod_{p \leq x}\left(1+\frac{2}{p}\right) .
$$

Now recall that $e^{t} \geq 1+t$ for every nonnegative $t$; this is clear from truncating the Taylor expansion $e^{t}=1+t+t^{2} / 2!+\ldots$. It follows that

$$
\prod_{p \leq x}\left(1+\frac{2}{p}\right) \leq \prod_{p \leq x} e^{2 / p}=\exp \left(\sum_{p \leq x} 2 / p\right) \leq \exp (2 C) .
$$

Consequently, the partial products

$$
\prod_{p \leq x}\left(1+\frac{1}{p}+\frac{1}{p^{2}}+\cdots\right)
$$

form a bounded, increasing sequence, which shows that we have condition (ii). We conclude that

$$
\sum_{n=1}^{\infty} \frac{1}{n}=\prod_{p} \frac{1}{1-\frac{1}{p}} \leq \exp (2 C)
$$

a contradiction.
Tweaking this argument, it is possible to derive an explicit lower bound on the partial sums $\sum_{p \leq x} 1 / p$ : Note that for $x \geq 2$,

$$
\begin{equation*}
\prod_{p \leq x} \frac{1}{1-\frac{1}{p}}=\sum_{n: p \mid n \Rightarrow p \leq x} \frac{1}{n} \geq \sum_{n \leq x} \frac{1}{n} \geq \log x \tag{1.4}
\end{equation*}
$$

From the upper bound $(1-1 / p)^{-1}=(1+1 /(p-1)) \leq \exp \left((p-1)^{-1}\right)$, we deduce (taking the logarithm of (1.4)) that $\sum_{p \leq x}(p-1)^{-1} \geq \log \log x$. To derive a lower bound for $\sum_{p \leq x} 1 / p$ from this, note that

$$
\begin{equation*}
\sum_{p \leq x} \frac{1}{p}=\sum_{p \leq x} \frac{1}{p-1}-\sum_{p \leq x}\left(\frac{1}{p-1}-\frac{1}{p}\right) \tag{1.5}
\end{equation*}
$$

$$
\geq \sum_{p \leq x} \frac{1}{p-1}-\sum_{n \geq 2}\left(\frac{1}{n-1}-\frac{1}{n}\right)=\left(\sum_{p \leq x} \frac{1}{p-1}\right)-1 \geq \log \log x-1
$$

The next two proofs also make use of the zeta function and its Euler factorization, but in a decidedly different manner.

Proof of J. Hacks. We need the well-known result, also due to Euler, that $\zeta(2)=\pi^{2} / 6$; a proof is sketched in Exercise 5 (for alternative arguments see [AZ04, Chapter 7], [Cha02]). Plugging $s=2$ into the Euler factorization (1.2) we obtain

$$
\frac{\pi^{2}}{6}=\zeta(2)=\prod_{p} \frac{1}{1-\frac{1}{p^{2}}}
$$

If there are only finitely many primes, then the product appearing here is a finite product of rational numbers, so that $\pi^{2} / 6$ must also be a rational number. But this is impossible, since $\pi$ is well known to be a transcendental number, i.e., not the root of any nonzero polynomial with rational coefficients. A weaker result, which suffices for the current argument, is the subject of Exercise 6 (cf. [AZ04, Chapter 6, Theorem 2]).

One can give a similar argument avoiding irrationality considerations:

Proof. We use not only that $\zeta(2)=\pi^{2} / 6$ but also that $\zeta(4)=\pi^{4} / 90$. (Again see Exercise 5.) Thus $\zeta(2)^{2} / \zeta(4)=5 / 2$. The Euler factorization (1.2) implies that

$$
\frac{5}{2}=\frac{\zeta(2)^{2}}{\zeta(4)}=\prod_{p}\left(1-p^{-4}\right)\left(1-p^{-2}\right)^{-2}=\prod_{p} \frac{p^{4}-1}{p^{4}} \frac{p^{4}}{\left(p^{2}-1\right)^{2}}=\prod_{p} \frac{p^{2}+1}{p^{2}-1},
$$

so that

$$
\frac{5}{2}=\frac{5}{3} \cdot \frac{10}{8} \cdot \frac{26}{24} \cdots
$$

If there are only finitely many primes, then the product on the right-hand side is a finite one and can be written as $M / N$, where $M=5 \cdot 10 \cdot 26 \cdots$ and $N=3 \cdot 8 \cdot 24 \cdots$. Then $M / N=5 / 2$, so $2 M=5 N$. Since $3 \mid N$, it must be that $3 \mid M$. But this cannot be: $M$ is a product of numbers of the form $k^{2}+1$, and no such number is a multiple of 3 .

Wagstaff has asked whether one can give a more elementary proof that $5 / 2=\prod_{p} \frac{p^{2}+1}{p^{2}-1}$. The discussion of this (open) question in [Guy04, B48] was the motivation for the preceding proof of Theorem 1.1.

## 5. Squarefree and smooth numbers

Recall that a natural number $n$ is said to be squarefree if it is not divisible by the square of any integer larger than 1 . The fundamental theorem of arithmetic shows that there is a bijection
\{finite subsets of the primes\} $\longleftrightarrow$ \{squarefree positive integers\},
given by sending

$$
S \longmapsto \prod_{p \in S} p
$$

So to prove the infinitude of the primes, it suffices to prove that there are infinitely many positive squarefree integers.
J. Perott's proof, 1881. We sieve out the non-squarefree integers from $1, \ldots, N$ by removing those divisible by $2^{2}$, then those divisible by $3^{2}$, etc. The number of removed integers is bounded above by

$$
\sum_{k=2}^{\infty}\left\lfloor N / k^{2}\right\rfloor \leq N \sum_{k=2}^{\infty} k^{-2}=N(\zeta(2)-1),
$$

so that the number of squarefree integers up to $N$, say $A(N)$, satisfies

$$
\begin{equation*}
A(N) \geq N-N(\zeta(2)-1)=N(2-\zeta(2)) \tag{1.6}
\end{equation*}
$$

At this point Perott uses the evaluation $\zeta(2)=\pi^{2} / 6$. However, it is simpler to proceed as follows: Since $t^{-2}$ is a decreasing function of $t$ on the positive real axis,

$$
\zeta(2)=1+\sum_{n=2}^{\infty} \frac{1}{n^{2}}<1+\sum_{n=1}^{\infty} \int_{n}^{n+1} \frac{d t}{t^{2}}=1+\int_{1}^{\infty} \frac{d t}{t^{2}}=2 .
$$

Referring back to (1.6), we see that $A(N) / N$ is bounded below by a positive constant. In particular, it must be that $A(N) \rightarrow \infty$ as $N \rightarrow \infty$.

Remark. As observed by Dressler [Dre75], Perott's argument also yields a lower bound on $\pi(N)$. Note that since every squarefree number $\leq N$ is a product of some subset of the $\pi(N)$ primes up to $N$, we have $2^{\pi(N)} \geq A(N)$. The argument above establishes that $A(N) \geq c N$ for $c=2-\zeta(2)>0$, and so $\pi(N) \geq \log N / \log 2+O(1)$.

For the next proof we need the following simple lemma:
Lemma 1.5. Every natural number $n$ can be written in the form $r s^{2}$, where $r$ and $s$ are natural numbers and $r$ is squarefree.

Proof. Choose the positive integer $s$ so that $s^{2}$ is the largest perfect square dividing $n$, and put $r=n / s^{2}$. We claim that $r$ is squarefree. Otherwise $p^{2} \mid r$ for some prime $p$. But then $(p s)^{2} \mid n$, contrary to the choice of $s$.

Erdős's proof of Theorem 1.1. Let $N$ be a positive integer. There are at most $\sqrt{N}$ squares not exceeding $N$ and at most $2^{\pi(N)}$ squarefree integers below this bound. So Lemma 1.5 implies that

$$
2^{\pi(N)} \sqrt{N} \geq N
$$

Dividing by $\sqrt{N}$ and taking logarithms yields the lower bound $\pi(N) \geq$ $\log N / \log 4$.

A modification of this argument leads to another proof that $\sum \frac{1}{p}$ diverges:

Erdős's proof of Theorem 1.4. Suppose that $\sum 1 / p$ converges. Then we can choose an $M$ for which

$$
\begin{equation*}
\sum_{p>M} \frac{1}{p}<\frac{1}{2} \tag{1.7}
\end{equation*}
$$

Keep this $M$ fixed.
Let $N$ be an arbitrary natural number. The estimate (1.7) implies that most integers up to $N$ factor completely over the primes not exceeding $M$.

Indeed, the number of integers not exceeding $N$ that have a prime factor $p>M$ is bounded above by

$$
\sum_{M<p \leq N}\left|\frac{N}{p}\right| \leq N \sum_{p>M} \frac{1}{p}<N / 2,
$$

so that more than $N / 2$ of the natural numbers not exceeding $N$ are divisible only by primes $p \leq M$.

We now show that there are too few integers divisible only by primes $p \leq M$ for this to be possible. There are at most $\sqrt{N}$ squares not exceeding $N$ and at most $C:=2^{\pi(M)}$ squarefree numbers composed only of primes not exceeding $M$. Thus there are at most $C \sqrt{N}$ natural numbers $\leq N$ having all their prime factors $\leq M$. But $C \sqrt{N}<N / 2$ once $N>4 C^{2}$.

In the last argument we needed an estimate for the number of integers up to a given point with only small prime factors. This motivates the following definition: Call a natural number $y$-smooth if all of its prime factors are bounded by $y$. We let $\Psi(x, y)$ denote the number of $y$-smooth numbers not exceeding $x$; i.e.,

$$
\begin{equation*}
\Psi(x, y):=\#\{n \leq x: p \mid n \Rightarrow p \leq y\} . \tag{1.8}
\end{equation*}
$$

Smooth numbers are important auxiliary tools in many number-theoretic investigations, and so there has been quite a bit of work on estimating the size of $\Psi(x, y)$ in various ranges of $x$ and $y$. (For a survey of both the applications and the estimates, see [Gra08b].) A trivial estimate yields an easy proof of Theorem 1.1.

Lemma 1.6. For $x \geq 1$ and $y \geq 2$, we have

$$
\Psi(x, y) \leq\left(1+\frac{\log x}{\log 2}\right)^{\pi(y)} .
$$

Proof. Let $k=\pi(y)$. By the fundamental theorem of arithmetic, $\Psi(x, y)$ is the number of $k$-tuples of nonnegative integers $e_{1}, \ldots, e_{k}$ with

$$
p_{1}^{e_{1}} p_{2}^{e_{2}} \cdots p_{k}^{e_{k}} \leq x
$$

This inequality requires $p_{i}^{e_{i}} \leq x$, so that

$$
e_{i} \leq \log x / \log p_{i} \leq \log x / \log 2,
$$

so that there are at most $1+\lfloor\log x / \log 2\rfloor$ possibilities for each $e_{i}$.
Since every positive integer not exceeding $N$ is a (possibly empty) product of primes not exceeding $N$,

$$
N=\Psi(N, N) \leq(1+\log N / \log 2)^{\pi(N)}
$$

It follows that

$$
\pi(N) \geq \frac{\log N}{\log (1+\log N / \log 2)}
$$

Taking some care to estimate the denominator, we obtain the lower bound

$$
\pi(N) \geq(1+o(1)) \frac{\log N}{\log \log N}
$$

which tends to infinity. Similar proofs of Theorem 1.1 have been given by Thue (1897), Auric (1915), Schnirelmann [Sch40, pp. 44-45], Chernoff [Che65], and Rubinstein [Rub93]. See also Exercise 17.

## 6. Sledgehammers!

In the spirit of the saying, "nothing is too simple to be made complicated," we finish off the first half of this chapter with two proofs of Theorem 1.1 that dip into the tool chest of higher mathematics.

The following "topological proof" is due to Furstenberg ([Fur55]):
Proof. We put a topology on $\mathbf{Z}$ by taking as a basis for the open sets all arithmetic progressions, infinite in both directions. (This is permissible since the intersection of two such progressions is either empty or is itself an arithmetic progression.) Then each arithmetic progression is both open and closed: it is open by choice of the basis, and it is closed since its complement is the union of the other arithmetic progressions with the same common difference. For each prime $p$, let $A_{p}=p \mathbf{Z}$, and define $A:=\bigcup_{p} A_{p}$. The set $\{-1,1\}=\mathbf{Z} \backslash A$ is not open. (Indeed, each open set is either empty or contains an arithmetic progression, so must be infinite.) It follows that $A$ is not closed. On the other hand, if there are only finitely many primes, then $A$ is a finite union of closed sets, and so it is closed.

Our next proof, due to L. Washington (and taken from [Rib96]) uses the machinery of commutative algebra. Recall that a Dedekind domain is an integral domain $R$ with the following three properties:
(i) $R$ is Noetherian: if $I_{1} \subset I_{2} \subset I_{3} \subset \cdots$ is an ascending chain of ideals of $R$, then there is an $n$ for which

$$
I_{n}=I_{n+1}=I_{n+2}=\cdots
$$

(ii) $R$ is integrally closed: if $K$ denotes the fraction field of $R$ and $\alpha \in K$ is the root of a monic polynomial with coefficients in $R$, then in fact $\alpha \in R$.
(iii) Every nonzero prime ideal of $R$ is a maximal ideal.

Proof. We use the theorem that a Dedekind domain with finitely many nonzero prime ideals is a principal ideal domain (see, e.g., [Lor96, Proposition III.2.12]) and thus also a unique factorization domain. The ring of integers $\mathfrak{O}_{K}$ of a number field $K$ is always a Dedekind domain; consequently, if $K$ does not possess unique factorization, then $\mathfrak{O}_{K}$ has infinitely many nonzero prime ideals. Each such prime ideal lies above a rational prime $p$, and for each prime $p$ there are at most $[K: \mathbf{Q}]$ prime ideals lying above it. It follows that there are infinitely many primes $p$, provided that there is a single number field $K$ for which $\mathfrak{O}_{K}$ does not possess unique factorization. And there is: If $K=\mathbf{Q}(\sqrt{-5})$, then

$$
6=(1+\sqrt{-5})(1-\sqrt{-5})
$$

is a well-known instance of the failure of unique factorization in $\mathfrak{O}_{K}=$ $\mathbf{Z}[\sqrt{-5}]$.

## 7. Prime-producing formulas

A mathematician is a conjurer who gives away his secrets. - J. H. Conway

Now that we know there are infinitely many primes, the next question is: Where are they hiding? Or, to ask a question that has ensnared many who have flirted with number theory: Is there a formula for producing primes? This line of inquiry, as natural as it seems, has not been very productive.

The following 1952 result of Sierpiński [Sie52] is representative of many in this subject. Let $p_{n}$ denote the $n$th prime number. Define a real number $\xi$ by putting

$$
\xi:=\sum_{n=1}^{\infty} p_{n} 10^{-2^{n}}=0.02030005000000070000000000000011 \ldots
$$

$\star$ Theorem 1.7. We have

$$
p_{n}=\left\lfloor 10^{2^{n}} \xi\right\rfloor-10^{2^{n-1}}\left\lfloor 10^{2 n-1} \xi\right\rfloor .
$$

This is, in the literal sense, a formula for primes. But while it may have some aesthetic merit, it must be considered a complete failure from the standpoint of utility; determining the number $\xi$ seems to require us to already know the sequence of primes. A similar criticism can be leveled against a result of Mills [Mil47], which asserts the existence of a real number $A>1$ with the property that $\left\lfloor A^{3^{n}}\right\rfloor$ is prime for each natural number $n$.

A more surprising way of generating primes was proposed by J. H. Conway [Con87]. Consider the following list of 14 fractions:

$$
\begin{array}{cccccccccccccc}
\text { A } & \text { B } & \text { C } & \text { D } & \text { E } & \text { F } & \text { G } & \text { H } & \text { I } & \text { J } & \text { K } & \text { L } & \text { M } & \text { N } \\
\frac{17}{91} & \frac{78}{85} & \frac{19}{51} & \frac{23}{38} & \frac{29}{33} & \frac{77}{29} & \frac{95}{23} & \frac{77}{19} & \frac{1}{17} & \frac{11}{13} & \frac{13}{11} & \frac{15}{2} & \frac{1}{7} & \frac{55}{1}
\end{array}
$$

Now run the following algorithm: Beginning with the number 2, look for the first (leftmost) fraction which can be multiplied by the current number to give an integer. Perform the multiplication and repeat. Whenever you reach a power of 2 , output the exponent. The first several (19) steps of the algorithm are

$$
\begin{gathered}
2 \mapsto 15 \mapsto 825 \mapsto 725 \mapsto 1925 \mapsto 2275 \mapsto 425 \mapsto 390 \mapsto 330 \mapsto 290 \mapsto 770 \\
\mapsto 910 \mapsto 170 \mapsto 156 \mapsto 132 \mapsto 116 \mapsto 308 \mapsto 364 \mapsto 68 \mapsto 4=2^{2}
\end{gathered}
$$

and so the first output is 2 . Fifty more steps yield

$$
2^{2} \mapsto 30 \mapsto 225 \mapsto 12375 \mapsto \cdots \mapsto 232 \mapsto 616 \mapsto 728 \mapsto 136 \mapsto 8=2^{3}
$$

and so the second output is 3 . After another 212 steps, we arrive at $32=2^{5}$, and so our third output is 5 .
$\star$ Theorem 1.8 (Conway). The sequence of outputs is exactly the sequence of primes in increasing order.

This is rather striking; the sequence of primes, which seems random in so many ways, is the output of a deterministic algorithm involving 14 fractions. But perhaps this should not come as such a shock. Most anyone who has experimented with programming knows that the primes are the output of a deterministic algorithm: Test the numbers $2,3,4, \ldots$ successively for primality, using (say) trial division for the individual tests. And actually, underneath the surface, this is exactly what is being done in Conway's algorithm. This sequence of 14 fractions encodes a simple computer program: The number $n$ is tested for divisibility first by $d=n-1$, then $d=n-2$, etc; as soon as a divisor is found, $n$ is incremented by 1 and the process is repeated. The game is rigged so that a power of 2 arises only when $d$ reaches 1, i.e., when $n$ is prime. Moreover, there is nothing special in Theorem 1.8 about the sequence of primes; an analogue of Theorem 1.8 can be proved for any recursive set. (Here a set of natural numbers $S$ is called recursive if there is an algorithm for determining whether a natural number belongs to $S$.) We conclude that while Conway's result is genuinely surprising, the surprise is that one can simulate computer programs with lists of fractions, and is in no way specific to the prime numbers.

## 8. Euler's prime-producing polynomial

The prime-producing functions we have been considering up to now have all been rather complicated. In some sense this is necessary; one can show that
any function which produces only primes cannot have too simple a form. We give only one early example of a result in this direction. (See [War30], [Rei43] for more theorems of this flavor.)

Theorem 1.9 (Goldbach). If $F(T) \in \mathbf{Z}[T]$ is a nonconstant polynomial with positive leading coefficient, then $F(n)$ is composite for infinitely many natural numbers $n$.

Proof. Suppose $F$ is nonconstant but that $F(n)$ is prime for all $n \geq N_{0}$, where $N_{0}$ is a natural number. Let $p=F\left(N_{0}\right)$; then $p$ divides $F\left(N_{0}+k p\right)$ for every positive integer $k$. But since $F$ has a positive leading coefficient, $F\left(N_{0}+k p\right)>p$ for every sufficiently large integer $k$, and so $F\left(N_{0}+k p\right)$ is composite, contrary to the choice of $N_{0}$.

Theorem 1.9 does not forbid the existence of polynomials $F$ which assume prime values over impressively long stretches. And indeed these do exist; a famous example is due to Euler, who observed that if $f(T)=$ $T^{2}+T+41$, then $f(n)$ is prime for all integers $0 \leq n<40$.

It turns out that Euler's observation, rather than being an isolated curiosity, is intimately connected with the theory of imaginary quadratic fields. We will prove the following theorem:

Theorem 1.10. Let $A \geq 2$, and set $D:=1-4 A$. Then the following are equivalent:
(i) $n^{2}+n+A$ is prime for all $0 \leq n<A-1$,
(ii) $n^{2}+n+A$ is prime for all $0 \leq n \leq \frac{1}{2} \sqrt{\frac{|D|}{3}}-\frac{1}{2}$,
(iii) the ring $\mathbf{Z}[(-1+\sqrt{D}) / 2]$ is a unique factorization domain.

The equivalence (i) $\Leftrightarrow$ (iii) is proved by Rabinowitsch in [Rab13], and is usually referred to as Rabinowitsch's theorem.

Remark. Since $n^{2}+n+A=(n+1 / 2)^{2}+(4 A-1) / 4$, (ii) can be rephrased as asserting that $(n+1 / 2)^{2}+|D| / 4$ is prime for every integer $n$ for which $|n+1 / 2| \leq \frac{1}{2} \sqrt{\frac{|D|}{3}}$. We will use this observation in the proof of Theorem 1.10.

Cognoscenti will recognize that $\mathbf{Z}[(-1+\sqrt{D}) / 2]$ is an order in the quadratic field $\mathbf{Q}(\sqrt{D})$. However, the proof of Theorem 1.10 presented here, due to Gyarmati (née Lanczi) [Lán65], [Gya83] and Zaupper [Zau83], requires neither the vocabulary of algebraic number theory nor the theory of ideals.

We begin the proof of Theorem 1.10 by observing that the bound on $n$ in (ii) is always at least as strict as the bound on $n$ in (i), which makes clear that (i) implies (ii). So it is enough to show that (ii) implies (iii)


Figure 1. The lattice $\mathbf{Z}+\mathbf{Z} \eta$ sitting inside $\mathbf{C}$. Here $A=2$ so that $D=-7$.
and that (iii) implies (i). To continue we need some preliminary results on the arithmetic of the rings $\mathbf{Z}[(-1+\sqrt{D}) / 2]$. These will be familiar to students of algebraic number theory, but we include full proofs for the sake of completeness.

Let $A \geq 2$ be an integer, and fix a complex root $\eta$ of $x^{2}+x+A$, so that (for an appropriate choice of the square root) $\eta=(-1+\sqrt{D}) / 2$. Since $\eta^{2}=-\eta-A$, it follows that

$$
\mathbf{Z}[\eta]=\mathbf{Z}+\mathbf{Z} \eta=\{x+y \eta: x, y \in \mathbf{Z}\} .
$$

For $\alpha \in \mathbf{Z}[\eta]$, we denote its complex conjugate by $\bar{\alpha}$. Observe that $\bar{\eta}=$ $-1-\eta$; consequently, $\mathbf{Z}[\eta]$ is closed under complex-conjugation. We define the norm of the element $\alpha=x+y \eta \in \mathbf{Z}[\eta]$ by

$$
\begin{aligned}
\mathcal{N}(\alpha): & =|\alpha|^{2} \\
& =\alpha \bar{\alpha}=x^{2}-x y+A y^{2} .
\end{aligned}
$$

Notice that the norm of $\alpha \in \mathbf{Z}[\eta]$ is always an integer and is positive whenever $\alpha \neq 0$. Moreover, since the complex absolute value is multiplicative, it is immediate that

$$
\mathcal{N}(\alpha \beta)=\mathcal{N}(\alpha) \cdot \mathcal{N}(\beta) \quad \text { for all } \quad \alpha, \beta \in \mathbf{Z}[\eta] .
$$

We now recall the requisite definitions from ring theory: If $\alpha, \beta \in \mathbf{Z}[\eta]$, we say that $\alpha$ divides $\beta$ if $\beta=\alpha \gamma$ for some $\gamma \in \mathbf{Z}[\eta]$. A nonzero element $\alpha \in \mathbf{Z}[\eta]$ is called a unit if $\alpha$ divides 1 . A nonunit element $\alpha \in \mathbf{Z}[\eta]$ is irreducible if whenever $\alpha=\beta \gamma$ with $\beta, \gamma \in \mathbf{Z}[\eta]$, then either $\beta$ is a unit or $\gamma$ is a unit. Finally, $\pi \in \mathbf{Z}[\eta]$ is called prime if whenever $\pi$ divides $\beta \gamma$ for $\beta, \gamma \in \mathbf{Z}[\eta]$, then either $\pi$ divides $\beta$ or $\pi$ divides $\gamma$.

Lemma 1.11. An element $\alpha \in \mathbf{Z}[\eta]$ is a unit precisely when $\mathcal{N}(\alpha)=1$. The only units in $\mathbf{Z}[\eta]$ are $\pm 1$.

Proof. If $\alpha$ is a unit, then $\mathcal{N}(\alpha) \cdot \mathcal{N}\left(\alpha^{-1}\right)=1$. Moreover, both $\mathcal{N}(\alpha)$ and $\mathcal{N}\left(\alpha^{-1}\right)$ are positive integers, so that $\mathcal{N}(\alpha)=\mathcal{N}\left(\alpha^{-1}\right)=1$. Conversely, if $\mathcal{N}(\alpha)=1$, then $\alpha \bar{\alpha}=1$, and so $\alpha$ is a unit. Finally, notice that if $y \neq 0$, then
$\mathcal{N}(x+y \eta)=x^{2}-x y+A y^{2}=(x-y / 2)^{2}+\frac{1}{4}(4 A-1) y^{2} \geq \frac{4 A-1}{4}>\frac{7}{4}>1$.
So $x+y \eta$ can be a unit only when $y=0$. In this case we must have $\mathcal{N}(x)=x^{2}=1$, and this occurs exactly when $x= \pm 1$.

Lemma 1.12. If $\alpha$ is a nonzero, nonunit element of $\mathbf{Z}[\eta]$, then $\alpha$ can be written as a product of irreducible elements of $\mathbf{Z}[\eta]$.

Proof. If the claim fails, there is a nonzero, nonunit $\alpha$ of smallest norm for which it fails. Clearly $\alpha$ is not irreducible, and so we can write $\alpha=\beta \gamma$, where $\beta$ and $\gamma$ are nonzero nonunits. Hence $\mathcal{N}(\alpha)=\mathcal{N}(\beta) \mathcal{N}(\gamma)$. Since $\mathcal{N}(\beta)$ and $\mathcal{N}(\gamma)$ are each larger than 1 , both $\mathcal{N}(\beta)$ and $\mathcal{N}(\gamma)$ must be smaller than $\mathcal{N}(\alpha)$. So by the choice of $\alpha$, both $\beta$ and $\gamma$ factor as products of irreducibles, and thus $\alpha$ does as well. This contradicts the choice of $\alpha$.

We can now prove one of the two outstanding implications:
Proof that (iii) $\Rightarrow$ (i). Let $\eta=(-1+\sqrt{D}) / 2$. Suppose $0 \leq n<A-1$. We have

$$
\begin{equation*}
n^{2}+n+A=(n-\eta)(n-\bar{\eta})=(n-\eta)(n+1+\eta) . \tag{1.9}
\end{equation*}
$$

Let $p$ be a prime dividing $n^{2}+n+A$. We claim that $p$ is not irreducible in $\mathbf{Z}[\eta]$. Indeed, since $\mathbf{Z}[\eta]$ is a unique factorization domain by hypothesis, if $p$ were irreducible, then $p$ would be prime. So from (1.9), we would have that $p$ divides $n-\eta$ or $n+1+\eta$. But this is impossible, since neither $n / p-\eta / p$ nor $(n+1) / p+\eta / p$ belongs to $\mathbf{Z}[\eta]=\mathbf{Z}+\mathbf{Z} \eta$.

Hence we can write $p=\alpha \beta$, where $\alpha, \beta \in \mathbf{Z}[\eta]$ and neither $\alpha$ nor $\beta$ is a unit. Taking norms, we deduce that $p^{2}=\mathcal{N}(p)=\mathcal{N}(\alpha) \mathcal{N}(\beta)$. Since $\alpha$ and $\beta$ are not units, we must have $\mathcal{N}(\alpha)=\mathcal{N}(\beta)=p$.

Write $\alpha=x+y \eta$ for integers $x, y$. Then $y \neq 0$ (since $p$ is a rational prime), and so

$$
p=N(\alpha)=x^{2}-x y+A y^{2}=(x-y / 2)^{2}+(A-1 / 4) y^{2} \geq A-1 / 4 .
$$

Thus (since $p$ is an integer) $p \geq A$. Moreover, since $0 \leq n<A-1$,

$$
n^{2}+n+A<(A-1)^{2}+(A-1)+A=(A-1) A+A=A^{2} .
$$

This shows that every prime divisor of $n^{2}+n+A$ exceeds its square root, so that $n^{2}+n+A$ is prime.

The proof of the remaining implication requires one more preliminary result:

Lemma 1.13. If $\pi$ is an element of $\mathbf{Z}[\eta]$ whose norm is a rational prime $p$, then $\pi$ is prime in $\mathbf{Z}[\eta]$.

Proof. We claim that $\mathbf{Z}[\eta] /(\pi)$ is isomorphic to $\mathbf{Z} / p \mathbf{Z}$. Since $\mathbf{Z} / p \mathbf{Z}$ is a field, this implies that $\pi$ generates a prime ideal of $\mathbf{Z}[\eta]$, which in turn implies that $\pi$ is prime. Let $\psi: \mathbf{Z} \rightarrow \mathbf{Z}[\eta] /(\pi)$ be the ring homomorphism defined by mapping $n$ to $n \bmod \pi$. Since $p=\pi \bar{\pi} \equiv 0(\bmod \pi)$, the kernel of $\psi$ contains the ideal $p \mathbf{Z}$. Since $p \mathbf{Z}$ is a maximal ideal, either $\psi$ is identically zero or the kernel of $\psi$ is precisely $p \mathbf{Z}$. Since $\pi$ is not a unit in $\mathbf{Z}[\eta], \psi(1)$ is nonzero, and so the kernel of $\psi$ is precisely $p \mathbf{Z}$. Hence $\mathbf{Z} / p \mathbf{Z}$ is isomorphic to the image of $\psi$. So the proof will be complete if we show that $\psi$ is surjective.

Write $\pi=r+s \eta$ for integers $r$ and $s$, and let $x+y \eta$ be an arbitrary element of $\mathbf{Z}[\eta]$. We can choose integers $a$ and $b$ for which

$$
m:=x+y \eta-\pi(a+b \eta) \in \mathbf{Z}
$$

Indeed, a short computation shows that this containment holds precisely when

$$
b(r-s)+a s=y,
$$

which is a solvable linear Diophantine equation in $a$ and $b$ since $\operatorname{gcd}(r-s, s)=$ $\operatorname{gcd}(r, s)=1$. Then $m \equiv x+y \eta(\bmod \pi)$, and so $\psi(m)=x+y \eta \bmod \pi$. Since $x+y \eta$ was arbitrary, $\psi$ is surjective as claimed.

Proof that (ii) $\Rightarrow$ (iii). Suppose that $n^{2}+n+A$ is prime for all

$$
0 \leq n \leq \frac{1}{2} \sqrt{\frac{|D|}{3}}-\frac{1}{2}
$$

We are to prove that $\mathbf{Z}[\eta]$ possesses unique factorization. Suppose otherwise, and let $\alpha$ be a nonzero, nonunit of minimal norm with two distinct factorizations into irreducibles, say

$$
\alpha=\pi_{1} \cdots \pi_{k}=\rho_{1} \cdots \rho_{j} .
$$

(Here distinct means that either $k \neq j$, or that $k=j$, but there is no way to reorder the $\pi_{i}$ so that each $\pi_{i}$ is a unit multiple of $\rho_{i}$.) By the minimality of $\mathcal{N}(\alpha)$, it is easy to see that none of the irreducibles in the first factorization can be a unit multiple of an irreducible in the second factorization. Consequently, none of the irreducibles appearing in either factorization can be prime in $\mathbf{Z}[\eta]$.


Figure 2. (Based on [Zau83].)

We can assume that $\mathcal{N}\left(\pi_{1}\right) \leq \mathcal{N}\left(\rho_{1}\right)$. (If this does not hold initially, interchange the two factorizations.) For $\xi, \gamma \in \mathbf{Z}[\eta]$ still to be chosen, define

$$
\begin{equation*}
\alpha^{\prime}:=\left(\rho_{1} \xi-\pi_{1} \gamma\right) \rho_{2} \cdots \rho_{j} . \tag{1.10}
\end{equation*}
$$

Then

$$
\begin{aligned}
\alpha^{\prime} & =\alpha \xi-\pi_{1} \frac{\alpha}{\rho_{1}} \gamma \\
& =\pi_{1}\left(\pi_{2} \cdots \pi_{k} \xi-\rho_{2} \cdots \rho_{j} \gamma\right) .
\end{aligned}
$$

Factoring the parenthetical expression, we deduce that $\alpha^{\prime}$ has a factorization into irreducibles where one of the irreducibles is $\pi_{1}$. We will choose $\xi$ and $\gamma$ so that $\pi_{1} \nmid \rho_{1} \xi$. Then $\pi_{1} \nmid \rho_{1} \xi-\pi_{1} \gamma$, and so we may deduce from (1.10) that $\alpha^{\prime}$ has a factorization into irreducibles, none of which is a unit multiple of $\pi_{1}$. So $\alpha^{\prime}$ possesses two distinct factorizations into irreducibles. If further, $\gamma$ and $\xi$ satisfy

$$
\mathcal{N}\left(\rho_{1} \xi-\pi_{1} \gamma\right)<\mathcal{N}\left(\rho_{1}\right),
$$

then $\mathcal{N}\left(\alpha^{\prime}\right)$ is smaller than $\mathcal{N}(\alpha)$, and so we have a contradiction to our choice of $\alpha$.

So it remains to show that it is possible to choose $\xi, \gamma \in \mathbf{Z}[\eta]$ with the following two properties:
(P1) $\pi_{1} \nmid \rho_{1} \xi$,
(P2) $\mathcal{N}\left(\rho_{1} \xi-\pi_{1} \gamma\right)<\mathcal{N}\left(\rho_{1}\right)$, or equivalently, $\left|\xi-\frac{\pi_{1}}{\rho_{1}} \gamma\right|<1$.
Since $\mathcal{N}\left(\pi_{1}\right) \leq \mathcal{N}\left(\rho_{1}\right)$, the complex number $\pi_{1} / \rho_{1}$ lies on or inside the unit circle. Suppose first that $\pi_{1} / \rho_{1}$ lies outside the shaded region indicated in


Figure 3. (Based on [Zau83].)
Figure 2. Then for either $\xi=1$ or $\xi=-1$, we have

$$
\left|\xi-\pi_{1} / \rho_{1}\right|<1
$$

Then (P1) and (P2) hold if we choose this value of $\xi$ and take $\gamma=1$. Note that $\pi_{1} \nmid \pm \rho_{1}$, since otherwise $\pi_{1}$ and $\rho_{1}$ would be unit multiples of each other, which we have already argued is not the case.

So we may assume that $\pi_{1} / \rho_{1}$ lies within the shaded region. Let $e_{1}$ be the ray from the origin making an angle of $60^{\circ}$ with the $x$-axis, and let $e_{2}$ be the ray from the origin making an angle of $120^{\circ}$ with that axis. Then the ray $e$ (say) from the origin through $\pi_{1} / \rho_{1}$ is contained within the $60^{\circ}$ angle determined by $e_{1}$ and $e_{2} .{ }^{2}$ Let $f$ be the horizontal line consisting of those complex numbers with imaginary part $\sqrt{|D|} / 2$; thus $f$ is the first horizontal line above the $x$-axis containing points of the lattice $\mathbf{Z}+\mathbf{Z} \eta$. Let $\mu$ be the

[^1]complex number corresponding to the intersection of $e$ and $f$. The angle determined by $e_{1}$ and $e_{2}$ cuts $f$ into a segment of length $\sqrt{|D| / 3}>1$, and so there is a point of $\mathbf{Z}+\mathbf{Z} \eta$ on $f$ within this angle. We choose such a point $\xi$ for which the distance from $\xi$ to $\mu$ is as small as possible. See Figure 3 .

We claim that the distance from $\xi$ to $e$ is strictly smaller than $\sqrt{3} / 2$. This is clear if both $\xi+1$ and $\xi-1$ fall within the angle determined by $e_{1}$ and $e_{2}$, since in that case, the distance from $\xi$ to $\mu$ must be at most $1 / 2$. So suppose that $\xi+1$ falls outside this angle; the case when $\xi-1$ falls outside is analogous. Then $\xi-1$ must lie within the given angle. Now, if $\xi$ is to the right of $\mu$, then in order that $\xi$ be at least as close to $\mu$ as $\xi-1$, it must be that the distance from $\xi$ to $\mu$ is at most $1 / 2$. So we can assume that $\xi$ falls to the left of $\mu$. This is the scenario depicted in Figure 3. In this case we we use the following argument: Let $\nu$ represent the intersection of $e_{1}$ and $f$; then the distance between $\xi$ and $\nu$ is smaller than 1 . Since $e_{1}$ makes an angle of $60^{\circ}$ with $f$, elementary trigonometry shows that the distance from $\xi$ to $e_{1}$ is strictly smaller than $\sqrt{3} / 2$. But the perpendicular line segment from $\xi$ to $e_{1}$ meets $e$. So the distance from $\xi$ to $e$ is also strictly smaller than $\sqrt{3} / 2$.

It follows that the unit disc centered at $\xi$ intersects $e$ in a segment of total length $>1$. (Indeed, let $\tau$ be the point on $e$ for which the line from $\xi$ to $\tau$ is perpendicular to $e$, so that the distance from $\xi$ to $\tau$ is strictly smaller than $\sqrt{3} / 2$. Then by the Pythagorean theorem, $\tau$ divides the segment in question into two parts, each of length $>1 / 2$.) Since $\left|\pi_{1} / \rho_{1}\right| \leq 1$, it follows that we can choose a rational integer $\gamma$ so that $\gamma \pi_{1} / \rho_{1}$ lies within the open unit disc centered at $\xi$.

We claim that with the above choices of $\xi$ and $\gamma$, both ( P 1 ) and ( P 2 ) hold. Condition (P2) is guaranteed by the choice of $\gamma$, so it remains only to verify ( P 1 ). For this it is enough to prove that $\xi$ is prime. Indeed, suppose that $\xi$ is prime but (P1) fails. Then

$$
\rho_{1} \xi=\pi_{1} \kappa
$$

for some $\kappa$. Since $\xi$ is prime, it must divide either $\pi_{1}$ or $\kappa$. But $\xi$ cannot divide $\pi_{1}$; if it did, then since $\pi_{1}$ is irreducible, we would have that $\pi_{1}$ is a unit multiple of $\xi$. But then $\pi_{1}$ would be prime since $\xi$ is prime. This contradicts the observation made above that none of the $\pi_{i}$ are prime. So $\xi$ must divide $\kappa$; but then dividing through by $\xi$ we find that $\pi_{1}$ divides $\rho_{1}$. That implies that $\pi_{1}$ and $\rho_{1}$ are unit multiples of each other, which again contradicts our initial observations.

Why should $\xi$ be prime? Since $\xi$ is a point of the lattice $\mathbf{Z}+\mathbf{Z} \eta$ lying on $f$, we have $\xi=n+\eta$ for some integer $n$. Moreover, since $\xi$ belongs to
the $60^{\circ}$ angle determined by $e_{1}$ and $e_{2}$, we find that

$$
|(n-1)+1 / 2|=|n-1 / 2| \leq \frac{1}{2} \sqrt{|D| / 3}
$$

But now (ii) of Theorem 1.10 implies that

$$
\begin{aligned}
\mathcal{N}(\xi) & =n^{2}-n+A \\
& =(n-1)^{2}+(n-1)+A
\end{aligned}
$$

is prime, so that $\xi$ is a prime element of $\mathbf{Z}[\eta]$ by Lemma 1.13.
A small amount of computation shows that condition (ii) of Theorem 1.10 holds for the values $A=2,3,5,11,17$, and 41 . This yields the following corollary:
Corollary 1.14. $\mathbf{Z}[(-1+\sqrt{D}) / 2]$ is a unique factorization domain for $D=$ $-7,-11,-19,-43,-67,-163$.

Checking larger values of $A$ does not appear to yield any more examples satisfying the conditions of Theorem 1.10. Whether or not the list in Corollary 1.14 is complete is known as the class number 1 problem; an equivalent question appears in Gauss's Disquisitiones (see [Gau86, Art. 303]). In 1933, Lehmer showed [Leh33] that any missing value of $A$ is necessarily large, in that $|D|>5 \cdot 10^{9}$. In 1934, Heilbronn \& Linfoot [HL34] showed that there is at most one missing value of $A$. Finally, in 1952, Heegner settled the problem, using new techniques from the theory of modular functions:

Theorem 1.15 (Heegner). If $A>41$, then $\mathbf{Z}[\eta]$ does not have unique factorization. Hence if $A \geq 2$ is an integer for which $n^{2}+n+A$ is prime for all $0 \leq n<A-1$, then $A \leq 41$.

For a modern account of Heegner's proof, see [Cox89, §12].

## 9. Primes represented by general polynomials

The result of the previous section leaves a very natural question unresolved: Does Euler's polynomial $T^{2}+T+41$, which does such a marvelous job of producing primes at the first several natural numbers $n$, represent infinitely many primes as $n$ ranges over the set of all positive integers? More generally, what can one say about the set of prime values assumed by a polynomial $F(T) \in \mathbf{Z}[T]$ ? In this section we survey the known results in this direction.
9.1. The linear case. Suppose first that $F(T)$ is linear, say $F(T)=a+$ $m T$, where $m>0$. Asking whether $F(n)$ is prime for infinitely many natural numbers $n$ amounts to asking whether the infinite arithmetic progression

$$
a+m, \quad a+2 m, \quad a+3 m, \quad a+4 m, \quad \ldots
$$

contains infinitely many primes - or, phrased in terms of congruences, whether or not there are infinitely many primes $p \equiv a(\bmod m)$.

This question is sometimes easy to answer. Let $d=\operatorname{gcd}(a, m)$. If $d>1$, then there are at most finitely many primes in the above progression, since every term is divisible by $d$, and so we have a negative answer to our query. So let us suppose that $d=1$. Then certain special cases can easily be settled in the affirmative. For example, if $a=-1$ and $m=4$, then we are asking for infinitely many primes $p \equiv-1(\bmod 4)$, and now we can mimic Euclid: If there are only finitely many such primes, say $p_{1}, \ldots, p_{k}$, form the number $N:=4 p_{1} \cdots p_{k}-1$. Since $N \equiv-1(\bmod 4)$, it must have at least one prime divisor $p \equiv-1(\bmod 4)$. But $p$ cannot be any of $p_{1}, \ldots, p_{k}$, and we have a contradiction. A similar argument works when $a=-1$ and $m=3$.

The general case of our problem is much more difficult. It turns out that whenever $\operatorname{gcd}(a, m)=1$, there are infinitely many primes $p \equiv a(\bmod m)$. This was proved by Dirichlet in 1837, by analytic methods. (One can view his argument as a far-reaching generalization of Euler's proof that the sum of the reciprocals of the primes diverges.) We will give a proof of Dirichlet's theorem in Chapter 4.

For now we content ourselves with some special cases of Dirichlet's theorem that follow from algebraic arguments. We noted above that an easy variant of Euclid's proof shows that there are infinitely many primes $p$ for which the residue class of $p$ avoids the trivial subgroup of the unit group $(\mathbf{Z} / 4 \mathbf{Z})^{\times}$, and similarly for $(\mathbf{Z} / 3 \mathbf{Z})^{\times}$. As observed by A. Granville (unpublished), we have the following general result:

Theorem 1.16. If $H$ is a proper subgroup of $(\mathbf{Z} / m \mathbf{Z})^{\times}$, then there are infinitely many primes $p$ for which $p \bmod m \notin H$.

Proof. Let $\mathcal{P}$ be the set of primes $p$ for which $p \bmod m \notin H$, and let $\mathcal{P}^{\prime}$ be the set of such primes not dividing $m$. Assuming $\mathcal{P}$ is finite, let $P$ be the product of the elements of $\mathcal{P}^{\prime}$. Fix an integer $a$ coprime to $m$ with $a \bmod m \notin H$ (which is possible since $H$ is a proper subgroup), and then choose a positive integer $n$ satisfying the congruences $n \equiv 1(\bmod P)$ and $n \equiv a(\bmod m)$. (Such a choice of $n$ is possible by the Chinese remainder theorem.) Since $n$ is coprime to $m P$, none of its prime divisors can come from $\mathcal{P}$, so that every prime $p$ dividing $n$ must be such that $p \bmod m \in H$. But since $H$ is closed under multiplication, this implies that $n \bmod m \in H$. This contradicts the choice of $a$.

If $F(T)$ is a nonzero polynomial with integer coefficients, we say that the prime $p$ is a prime divisor of $F$ if $p$ divides $F(n)$ for some integer $n$. The following useful lemma is due to Schur [Sch12]:

Lemma 1.17. Let $F(T)$ be a nonconstant polynomial with integer coefficients. Then $F$ has infinitely many prime divisors.

Proof. If $F(0)=0$, then every prime is a prime divisor of $F$. So we can assume that the constant term $c_{0}$ (say) of $F(T)$ is nonzero. Then $F\left(c_{0} T\right)=c_{0} G(T)$ for some nonconstant polynomial $G(T)$ with constant term 1. It is enough to show that $G$ has infinitely many prime divisors. Suppose that $p_{1}, \ldots, p_{k}$ is a list of prime divisors of $G$. For $m$ sufficiently large, we have $\left|G\left(m p_{1} \cdots p_{k}\right)\right|>1$, so that there must be some prime $p$ dividing $G\left(m p_{1} \cdots p_{k}\right)$. Then $p$ is a prime divisor of $G$ and $p$ is not equal to any of the $p_{i}$, since $G\left(m p_{1} \cdots p_{k}\right) \equiv 1\left(\bmod p_{i}\right)$ for each $1 \leq i \leq k$. So no finite list of prime divisors of $G$ can be complete.

For example, let $F(T)=T^{2}+1$. If $p$ divides $n^{2}+1$, then $n^{2} \equiv$ $-1(\bmod p)$, and so either $p=2$ or $p \equiv 1(\bmod 4)$. So Lemma 1.17 implies that there are infinitely many primes $p \equiv 1(\bmod 4)$. Similarly, if $F(T)=T^{2}+T+1$, then any prime divisor $p$ of $F$ is such that $p \equiv 1(\bmod 3)$, and so there are infinitely many primes $p \equiv 1(\bmod 3)$. Combining this with our earlier results, we have proved Dirichlet's theorem for all progressions modulo 3 and modulo 4.

These examples are special cases of the following construction: Recall that the mth cyclotomic polynomial is defined by

$$
\Phi_{m}(T)=\prod_{\substack{1 \leq k \leq m \\ \operatorname{gcd}(k, m)=1}}\left(T-e^{2 \pi i k / m}\right)
$$

i.e., $\Phi_{m}(T)$ is the monic polynomial in $\mathbf{C}[T]$ whose roots are precisely the primitive $m$ th roots of unity, each occurring with multiplicity 1 . For example, $\Phi_{4}(T)=T^{2}+1$ and $\Phi_{3}(T)=T^{2}+T+1$.

We will apply Lemma 1.17 to $\Phi_{m}$ to deduce that there are infinitely many primes $p \equiv 1(\bmod m)$. To apply Lemma 1.17 , we need that the coefficients of $\Phi_{m}(T)$ are not merely complex numbers, but in fact integers.

Lemma 1.18. For each positive integer $m$, the polynomial $\Phi_{m}(T)$ has integer coefficients.

Proof. For each $m$ we have the factorization

$$
\begin{equation*}
T^{m}-1=\prod_{d \mid m} \Phi_{d}(T) \tag{1.11}
\end{equation*}
$$

To see this, note that $T^{m}-1=\prod_{\zeta^{m}=1}(T-\zeta)$. Since the set of $m$ th roots of unity is the disjoint union of the primitive $d$ th roots of unity, taken over
those $d$ dividing $m$, we have (1.11). Applying Möbius inversion to (1.11) yields

$$
\Phi_{m}(T)=\prod_{d \mid m}\left(T^{d}-1\right)^{\mu(m / d)}=\frac{\prod_{d \mid m, \mu(m / d)=1}\left(T^{d}-1\right)}{\prod_{d \mid m, \mu(m / d)=-1}\left(T^{d}-1\right)}=\frac{F}{G},
$$

say. Now $F$ and $G$ are monic polynomials in $\mathbf{Z}[T]$ with $G \neq 0$, and so we can write

$$
\begin{equation*}
F=G Q+R, \tag{1.12}
\end{equation*}
$$

where $Q, R \in \mathbf{Z}[T]$ and $\operatorname{deg} R<\operatorname{deg} Q$. Of course (1.12) remains valid over $\mathbf{C}[T]$ and expresses in that ring one result of division by $G$. But we know that over $\mathbf{C}[T]$, we have $F=G \Phi_{m}$, so that $G$ goes into $F$ with no remainder. By the uniqueness of quotient and remainder in the division algorithm for polynomials, we must have $R=0$ above. Consequently, $\Phi_{m}=F / G=Q \in$ $\mathbf{Z}[T]$.
Lemma 1.19. If $p$ is a prime divisor of $\Phi_{m}$, then either $p \mid m$ or $p \equiv$ $1(\bmod m)$.

Proof. If $p$ is a prime divisor of $\Phi_{m}$, then $p$ divides $\Phi_{m}(n)$ for some integer $n$. Since the cyclotomic polynomials have integer coefficients, it follows from (1.11) that $p \mid \prod_{d \mid m} \Phi_{d}(n)=n^{m}-1$, so that the order of $n$ modulo $p$ is a divisor of $m$.

Suppose now that $p$ does not divide $m$. We claim that in this case, $m$ is the precise order of $n$ modulo $p$. Thus $m$ divides $p-1$, whence $p \equiv$ $1(\bmod m)$. To prove the claim, suppose for the sake of contradiction that $f<m$ is the exact order of $n \bmod p$. Then $f$ is a proper divisor of $m$. Moreover, $p$ divides $n^{f}-1=\prod_{e \mid f} \Phi_{e}(n)$, so that $p$ divides $\Phi_{e}(n)$ for some $e \mid$ $f$. Hence the residue class $n \bmod p$ is a zero of both $\Phi_{e}(T)$ and $\Phi_{m}(T)$. The polynomials $\Phi_{e}$ and $\Phi_{m}$ both appear in the factorization (1.11) of $T^{m}-1$, so that $T^{m}-1$ has a zero of order $\geq 2$ over $\mathbf{Z} / p \mathbf{Z}$. But $T^{m}-1$ has no multiple roots over $\mathbf{Z} / p \mathbf{Z}$, since $T^{m}-1$ has no roots in common with its derivative $m T^{m-1}$.

Since only finitely many primes divide $m$, Lemmas 1.17 and 1.19 have the following corollary:

Corollary 1.20. For each natural number $m$, there are infinitely many primes $p \equiv 1(\bmod m)$.

This proof of Corollary 1.20 is essentially due to Wendt [Wen95].
How far can one take this algebraic approach? The following result is due to Schur (op. cit.).
$\star$ Theorem 1.21. Let $m$ be a positive integer and let $H$ be a subgroup of $(\mathbf{Z} / m \mathbf{Z})^{\times}$. There is a nonconstant polynomial $F(T) \in \mathbf{Z}[T]$ with the following property: Every prime divisor p of $F$, with finitely many exceptions, satisfies $p \bmod m \in H$. Consequently, there are infinitely many primes $p$ for which $p \bmod m \in H$.

When $H$ is the trivial subgroup we have just seen that $F:=\Phi_{m}$ satisfies the conclusion of Theorem 1.21.

Schur gave an elementary proof of Theorem 1.21 requiring only familiarity with the theory of finite fields. A less elementary proof is outlined in Exercise 20. When $m$ is a prime number, Theorem 1.21 is contained in the results of Chapter 2 (see, in particular, Theorem 2.15).

Suppose that $a$ and $m$ satisfy $a^{2} \equiv 1(\bmod m)$, where $a \not \equiv 1(\bmod m)$. Applying Theorem 1.21 to the 2-element subgroup of $(\mathbf{Z} / m \mathbf{Z})^{\times}$generated by $a \bmod m$, we obtain a polynomial $F(T)$ all of whose prime divisors (with finitely many exceptions) satisfy either $p \equiv 1(\bmod m)$ or $p \equiv a(\bmod m)$. Schur showed (op. cit.) that if there is a single, suitably large prime $p \equiv$ $a(\bmod m)$, then the polynomial $F$ he constructs cannot have all (or even all but finitely many) of its prime divisors from the progression $1 \bmod m$. (See the first example below for an illustration of how this works.) So $F$ must have infinitely many prime divisors $p \equiv a(\bmod m)$.

Since Dirichlet's theorem is true, there is always a suitably large prime $p \equiv a(\bmod m)$ to be used in Schur's argument, and so in principle, it is possible to give a purely algebraic proof of Dirichlet's theorem for any progression $a \bmod m$ satisfying $a^{2} \equiv 1(\bmod m)$. Moreover, this is best possible in the following sense:
$\star$ Theorem 1.22 (Murty [Mur88, MT06]). Suppose $m$ is a positive integer. If $F$ is a nonconstant polynomial with the property that every prime divisor $p$ of $F$, with finitely many exceptions, satisfies $p \equiv 1(\bmod m)$ or $p \equiv a(\bmod m)$, then $a^{2} \equiv 1(\bmod m)$.

The proof of Theorem 1.22 rests on rather deep results in algebroanalytic number theory. The principal tool required is the Chebotarev density theorem, which is a far-reaching generalization of Dirichlet's theorem. See [SL96] for a down-to-earth discussion of Chebotarev's result.

Example. As an easy example of Schur's method, consider the problem of showing that there are infinitely many primes $p \equiv 3(\bmod 8)$. We start by taking $F(T):=T^{2}+2$. From the elementary theory of quadratic residues we have that each odd prime divisor of $F(T)$ satisfies $p \equiv 1$ or $3(\bmod 8)$. Now we observe that there is at least one prime in the residue class $3(\bmod 8)$,
namely 11. We replace $T$ by $4 T+3$ and so obtain from $F$ the polynomial

$$
G(T)=F(4 T+3)=16 T^{2}+24 T+11=8\left(2 T^{2}+3 T\right)+11 .
$$

Then every prime divisor of $G$ belongs to either the residue class $1 \bmod 8$ or $3 \bmod 8$. Moreover, for each positive integer $n$, there is at least one prime $p \equiv 3(\bmod 8)$ for which $p \mid G(n)$, since $G(n) \equiv 3(\bmod 8)$. We will show that $G$ (and hence also $F$ ) must have infinitely many prime divisors from the residue class $3 \bmod 8$. Suppose otherwise, and let $p_{1}, p_{2}, \ldots, p_{k}$ be a complete list of the prime divisors $p \equiv 3(\bmod 8)$ of $G$. For each $p_{i}$, choose an integer $n_{i}$ for which $G\left(n_{i}\right) \not \equiv 0\left(\bmod p_{i}\right)$. (This is possible since $G$ has at most two roots modulo $p_{i}$.) If $n$ is a positive integer chosen by the Chinese remainder theorem to satisfy $n \equiv n_{i}\left(\bmod p_{i}\right)$ for all $1 \leq i \leq k$, then $G(n)$ cannot be divisible by any of $p_{1}, \ldots, p_{k}$. So $G(n)$ must have a prime divisor from the residue class $3 \bmod 8$ other than $p_{1}, \ldots, p_{k}$, a contradiction.

Example. Since every integer $a$ coprime to 24 satisfies $a^{2} \equiv 1(\bmod 24)$, it is in principle possible to give an algebraic proof of Dirichlet's theorem for progressions with common difference 24 . The details in this case have been completely worked out by Bateman \& Low [BL65]. We leave to the reader the task of showing that 24 is the largest modulus $m$ with the property that $a^{2} \equiv 1(\bmod m)$ for each $a$ coprime to $m$.

### 9.2. Hypothesis H.

I do not mean to deny that there are mathematical truths, morally certain, which defy and will probably to the end of time continue to defy proof, as, e.g., that every indecomposable polynomial function must represent an infinitude of primes. - J. J. Sylvester [Syl88]
There are two natural directions we might head in if we hope to generalize Dirichlet's result: First, we might inquire about simultaneous prime values of several linear polynomials. One has to be careful here, of course. For example, we cannot hope that there are infinitely many $n$ for which both $n$ and $n+1$ are prime, because one of these two numbers is always even! However, if instead of $n$ and $n+1$ we consider $n$ and $n+2$, then this obstruction disappears, and we arrive at the following famous conjecture:

Conjecture 1.23 (Twin prime conjecture). There are infinitely many natural numbers $n$ for which both $n$ and $n+2$ are prime.

Alternatively, we might accept the restriction of working with a single polynomial, but hope to treat polynomials of higher degree. The following conjecture of Euler, which appears in correspondence with Goldbach, fits nicely into this framework:

Conjecture 1.24 (Euler). There are infinitely many natural numbers $n$ for which $n^{2}+1$ is prime.

Similarly, it seems reasonable to conjecture that our old friend, $T^{2}+T+$ 41, represents infinitely many primes. Once again, formulating conjectures of this type requires some care; if $n^{2}+1$ or $n^{2}+n+41$ is replaced by $n^{2}+n+2$, then the statement corresponding to Euler's conjecture is false, since $n^{2}+n+2$ is always even.

Suppose more generally that $F_{1}(T), \ldots, F_{r}(T) \in \mathbf{Z}[T]$ are nonconstant polynomials, each with positive leading coefficient. We can ask when it is the case that $F_{1}(n), \ldots, F_{r}(n)$ are simultaneously prime for infinitely many natural numbers $n$. Evidently if this is to be the case, then we must suppose that each $F_{i}$ is irreducible over $\mathbf{Z}$. The example of $r=2$ and $F_{1}(T)=T$, $F_{2}(T)=T+1$ shows that this is not sufficient, as does the example of $r=1$ and $F_{1}(T)=T^{2}+T+2$. What goes wrong in these examples is that there is a local obstruction: If we put $G(T):=\prod_{i=1}^{r} F_{i}(T)$, then $G(n)$ is always even. In 1958, Schinzel conjectured (see [SS58]) that these are the only remaining obstructions to be accounted for:

Conjecture 1.25 (Schinzel's "Hypothesis H"). Suppose $F_{1}(T), \ldots, F_{r}(T) \in$ $\mathbf{Z}[T]$ are nonconstant and irreducible and that each $F_{i}$ has a positive leading coefficient. Put $G(T):=\prod_{i=1}^{r} F_{i}(T)$, and suppose that there is no prime $p$ which divides $G(n)$ for every integer $n$. Then $F_{1}(n), F_{2}(n), \ldots, F_{r}(n)$ are simultaneously prime for infinitely many natural numbers $n$.

The hypothesis on $G$ is necessary: Suppose that $p$ is a (fixed) prime which divides $G(n)$ for each $n$. Then $p$ divides some $F_{i}(n)$ for each $n$. But for large $n$, each $F_{i}(n)>p$, and so for large $n$, some $F_{i}(n)$ is composite.

The twin prime conjecture corresponds to choosing $r=2, F_{1}(T)=T$, and $F_{2}(T)=T+2$ in Hypothesis H. Taking instead $r=1$ and $F_{1}(T)=$ $T^{2}+1$, we recover Euler's Conjecture 1.24. Despite substantial attention, both the twin prime conjecture and Conjecture 1.24 remain open. Even more depressing, no case of Hypothesis H has ever been shown to hold except when $r=1$ and $F_{1}(T)$ is linear, when Hypothesis H reduces to Dirichlet's theorem!

Sieve methods, which we introduce in Chapter 6, can be used to obtain certain approximations to Hypothesis H. We give two examples: A theorem of Chen [Che73] asserts that there are infinitely many primes $p$ for which $p+2$ is either prime or the product of two primes. And Iwaniec [Iwa78] has shown that there are infinitely many $n$ for which $n^{2}+1$ is either prime or the product of two primes. (This latter result applies also to $n^{2}+n+41$, and in fact to any quadratic obeying the conditions of Hypothesis H.)

## 10. Primes and composites in other sequences

We conclude by discussing the occurrence of primes in other sequences of interest. Results in this area are rather thin on the ground, and so we content ourselves with a smattering of problems and results meant to showcase our collective ignorance.

One sequence that has received much attention is that of the Mersenne numbers $2^{n}-1$. The occurrence of primes in this sequence has long been of interest in view of Euclid's result that if $2^{n}-1$ is prime, then $2^{n-1}\left(2^{n}-1\right)$ is a perfect number. (Here a number is called perfect if it is the sum of its proper divisors.) Since $2^{d}-1$ divides $2^{n}-1$ whenever $d$ divides $n$, for $2^{n}-1$ to be prime it is necessary that $n$ be prime. At first glance it appears that $2^{p}-1$ is often prime; 7 of the first 10 primes $p$ have this property. However, the tide quickly turns: Of the 78498 primes $p$ up to $10^{6}$, only 31 yield primes. As of February 2009, there are 46 known primes of the form $2^{p}-1$, the largest corresponding to $p=43112609$. It is not clear from this data whether or not we should expect infinitely many primes of this form, but probabilistic considerations to be discussed in Chapter 3 suggest that we should:

Conjecture 1.26. For infinitely many primes $p$, the number $2^{p}-1$ is prime.
Unfortunately, this conjecture seems far beyond reach. In fact, we know disturbingly little about the numbers $2^{p}-1$; perhaps the most striking illustration of this is that even the following modest conjecture remains unproved:

Conjecture 1.27. For infinitely many primes $p$, the number $2^{p}-1$ is composite.

We may also change the " - " sign to a " + " and consider primes of the form $2^{n}+1$. Since $2^{d}+1$ divides $2^{n}+1$ when $n / d$ is odd, we see that $2^{n}+1$ can be prime only if $n$ is a power of 2 . This leads us to consider the Fermat numbers $F_{m}=2^{2^{m}}+1$. The attentive reader will recall that these numbers appeared already in Goldbach's proof of Theorem 1.1. For $m=0,1,2,3$, and 4 , the numbers $F_{m}$ are prime:
$2^{2^{0}}+1=3, \quad 2^{2^{1}}+1=5, \quad 2^{2^{2}}+1=17, \quad 2^{2^{3}}+1=257, \quad 2^{2^{4}}+1=65537$.
Fermat was intuitively certain that $F_{m}$ is prime for all $m \geq 0$, and expressed this belief in letters to his contemporaries; but in 1732 Euler discovered the factorization

$$
2^{2^{5}}+1=641 \cdot 6700417 .
$$

It is now known that $F_{m}$ is composite for $5 \leq m \leq 32$, and (for the same probabilistic reasons alluded to above) it is widely believed that $F_{m}$ is composite for every $m \geq 5$. So much for intuition! Despite this widespread belief, the following conjecture appears intractable:
Conjecture 1.28. The Fermat number $F_{m}$ is composite for infinitely many natural numbers $m$.

Similarly, for each even natural number $a$, one can look for primes in the sequence $a^{2^{m}}+1$. Again we believe that there should be at most finitely many, but again the analogue of Conjecture 1.28 seems impossibly difficult! Indeed, there is no specific even number $a$ for which we can prove that $a^{2^{m}}+1$ is composite infinitely often. This is a somewhat odd state of affairs in view of the following amusing theorem of Schinzel [Sch63]:
Theorem 1.29. Suppose that infinitely many of the Fermat numbers $F_{j}$ are prime. If $a>1$ is an integer not of the form $2^{2^{r}}$ (where $r \geq 0$ ), then $a^{2^{j n}}+1$ is composite for infinitely many natural numbers $m$.

Proof. Fix an integer $a>1$ not of the form $2^{2^{r}}$. Let $M_{0}$ be an arbitrary positive integer. We will show that $a^{2^{m}}+1$ is composite for some $m \geq M_{0}$.

Let $F_{j}$ be a prime Fermat number not dividing $a\left(a^{2^{M_{0}}}-1\right)$. Since $a$ is coprime to $F_{j}$, Fermat's little theorem implies that

$$
a^{F_{j}-1}=a^{2^{2^{j}}} \equiv 1\left(\bmod F_{j}\right) .
$$

Since $F_{j} \nmid a^{2^{M_{0}}}-1$, we must have $M_{0}<2^{j}$. So we can write

$$
\begin{aligned}
a^{F_{j}-1}-1 & =a^{2^{2^{j}}}-1 \\
& =\left(a^{2^{M_{0}}}-1\right)\left(a^{2^{M_{0}}}+1\right)\left(a^{2^{M_{0}+1}}+1\right)\left(a^{2^{M_{0}+2}}+1\right) \cdots\left(a^{2^{2^{j}-1}}+1\right) .
\end{aligned}
$$

Since $F_{j}$ divides $a^{F_{j}-1}-1$ but not $a^{2^{M_{0}}}-1$, it must be that $F_{j}$ divides $a^{2^{m}}+1$ for some $M_{0} \leq m<2^{j}$. We cannot have $a^{2^{m}}+1=F_{j}$, since $a$ is not of the form $2^{2^{r}}$, and so $a^{2^{m}}+1$ is composite.

In connection with Fermat-type numbers the following result of Shapiro \& Sparer [SS72] merits attention (cf. [Sha83, Theorem 5.1.5]). It shows (in particular) that the doubly exponential sequences $a^{2^{m}}+1$ are unusually difficult to handle among sequences of the same general shape:
$\star$ Theorem 1.30. Suppose $a, b$, and $c$ are integers, and that $a, b>1$. If $c$ is odd, then

$$
a^{b^{m}}+c
$$

is composite for infinitely many $m \in \mathbf{N}$, except possibly in the case when a is even, $c=1$, and $b=2^{k}$ for some $k \geq 1$. If $c$ is even, there are infinitely many such $m$ except possibly when $a$ is odd and $c=2$.

The reader should note that the Shapiro-Sparer paper contains several other attractive results on composite numbers in various sequences.

We close this section by considering the sequence of shifted factorials $n!+1$. Here we can easily obtain infinitely many composite terms, since Wilson's theorem implies that $(p-1)!+1$ is composite for each $p>3$. The following pretty theorem of Schinzel [ $\mathbf{S c h} \mathbf{6 2 b}$ ] generalizes this result:
Theorem 1.31. Let $\alpha$ be a positive rational number. Then there are infinitely many $n$ for which $\alpha \cdot n!+1$ is composite.
Lemma 1.32. Let $p$ be a prime and let $r$ and $s$ be positive integers. Then for $0 \leq i \leq p-1$, we have

$$
p\left|s i!+(-1)^{i+1} r \Longleftrightarrow p\right| r(p-1-i)!+s .
$$

Proof. By Wilson's theorem,

$$
\begin{aligned}
-1 \equiv(p-1)! & =(p-1)(p-2) \cdots(p-i)(p-i-1)! \\
& \equiv(-1)^{i} i!(p-1-i)!(\bmod p),
\end{aligned}
$$

so that $(p-1-i)!i!\equiv(-1)^{i+1}(\bmod p)$. Since $p$ and $(p-1-i)!$ are coprime,

$$
\begin{aligned}
p \mid s i!+(-1)^{i+1} r & \Longleftrightarrow p \mid s(p-1-i)!i!+(-1)^{i+1} r(p-1-i)! \\
& \Longleftrightarrow p \mid(-1)^{i+1} s+(-1)^{i+1} r(p-1-i)! \\
& \Longleftrightarrow p \mid s+r(p-1-i)!.
\end{aligned}
$$

Proof of Theorem 1.31. Write $\alpha=r / s$, where $r$ and $s$ are relatively prime positive integers. Assume $l \in \mathbf{N}$ and $l \geq r / 2$. Then ( $4 l$ )! $\alpha^{-1}$ is an integer divisible by both 4 and $r$. Since $4 \mid(4 l)!\alpha^{-1}$, we can choose a prime $p_{l} \equiv-1(\bmod 4)$ with

$$
p_{l} \mid(4 l)!\alpha^{-1}-1 .
$$

Because $r \mid(4 l)!\alpha^{-1}$, necessarily $p_{l} \nmid r$. Since

$$
\begin{equation*}
p_{l} \mid r\left((4 l)!\alpha^{-1}-1\right)=s(4 l)!-r, \tag{1.13}
\end{equation*}
$$

we must have $p_{l}>4 l$. From Lemma 1.32 (with $i=4 l$ ) and (1.13), we find that

$$
\begin{equation*}
p_{l} \mid r\left(p_{l}-4 l-1\right)!+s \tag{1.14}
\end{equation*}
$$

Since $p_{l} \nmid r$, (1.14) implies that $p_{l} \nmid s$, and so

$$
p_{l} \mid N_{l}:=\alpha\left(p_{l}-4 l-1\right)!+1
$$

whenever $N_{l}$ is an integer. This happens for all large $l$ : Indeed, from (1.14) we have $N_{l} \geq p_{l} / s \geq 4 l / s$, so that $N_{l} \rightarrow \infty$ with $l$, which is only possible if $p_{l}-4 l-1 \rightarrow \infty$ with $l$. But $N_{l}$ is an integer whenever $p_{l}-4 l-1 \geq s$.

Finally, notice that for large $l$, we cannot have $p_{l}=N_{l}$, since $p_{l} \equiv$ $-1(\bmod 4)$ while $N_{l} \equiv 1(\bmod 4)$. Thus $N_{l}$ is a composite integer of the
form $\alpha \cdot n!+1$. Letting $l \rightarrow \infty$, we obtain infinitely many composite numbers of this form.

## Notes

Most of the proofs discussed for the infinitude of the primes may be found in [Dic66, Chapter XVIII] or [Nar00, §1.1]. For other compilations, see [Rib96, Chapter 1], [FR07, Chapter 3], and [Moh79]. An amusing version of Euclid's proof, couched in the language of nonstandard analysis, is presented in [Gol98, pp. 57-58]. Additional elementary proofs of the stronger result that $\sum 1 / p$ diverges may be found in [Bel43], [Mos58], and the survey [VE80].

The following result of Matijasevich and Putnam provides an interesting contrast to Goldbach's theorem (Theorem 1.9): There is a polynomial with integral coefficients such that the set of primes coincides with the set of positive values assumed by this polynomial, as the variables range over the nonnegative integers. (An explicit example of such a polynomial, in 26 variables, was produced by Jones et al. [JSWW76].) Yet upon inspection we realize we are once again looking at a result that properly belongs not to number theory but to computability theory (or logic); an analogous statement is true if we replace the set of primes with any listable set. Here a set of positive integers $S$ is called listable if there is a computer program which, when left running forever, outputs precisely the elements of $S$. A very approachable introduction to this circle of ideas is Matijasevich's article [Mat99]; for complete details see [Mat93].

In connection with the results of $\S 8$, we cannot resist pointing out the remarkable identity

$$
e^{\pi \sqrt{163}}=262537412640768743.99999999999925 \ldots,
$$

which shows that $e^{\pi \sqrt{163}}$ is very nearly an integer. We sketch the explanation, which comes from the theory of modular functions; for details one may consult [Cox89, §11]. Every lattice $L \subset \mathbf{C}$ has a so-called $j$-invariant $j(L)$, and $j\left(L_{1}\right)=j\left(L_{2}\right)$ precisely when $L_{1}$ and $L_{2}$ are homothetic, i.e., when one can be obtained from the other by rotation and scaling. We view $j$ as a function on the upper half-plane $\{z \in \mathbf{C}: \Im(z)>0\}$ by defining $j(\tau)$ as $j(L)$, where $L$ is the lattice spanned by 1 and $\tau$. It turns out that $j$ is then holomorphic on the upper half-plane. Moreover, since 1 and $\tau$ determine the same lattice as 1 and $\tau+1$, we have $j(\tau)=j(\tau+1)$. This shows that $j(\tau)$ is holomorphic as a function of $q=e^{2 \pi i \tau}$ in the punctured disc $0<|q|<1$, and so $j$ has a Laurent expansion. It turns out that this expansion starts

$$
j(\tau)=\frac{1}{q}+744+196884 q+\cdots
$$

so that $j(\tau) \approx 1 / q+744$ for small $q$. Now for the coup de grâce: One can show that if $K$ is an imaginary quadratic field with integral basis $1, \tau$, then $j(\tau)$ is an algebraic integer of degree exactly $h(K)$, the class number of $K$. In particular, if $K$ has class number 1 , then $j(\tau)$ is a rational integer. The main theorem of $\S 8$ implies that $K=\mathbf{Q}(\sqrt{-163})$ has class number 1 , and so $j(\tau) \in \mathbf{Z}$ for $\tau=\frac{1+i \sqrt{163}}{2}$. This value of $\tau$ corresponds to $q=$ $-1 / \exp (\pi \sqrt{163})$, so that

$$
e^{\pi \sqrt{163}} \approx j(\tau)-744 \in \mathbf{Z}
$$

We remark that $e^{\pi \sqrt{163}}$ is actually transcendental, as may be deduced from the following theorem of Gelfond and Schneider (noting that $e^{\pi \sqrt{163}}=$ $\left.(-1)^{i \sqrt{163}}\right)$ : If $\alpha$ and $\beta$ are algebraic numbers, where $\alpha \neq 0$ and $\beta$ is irrational, then $\alpha^{\beta}$ is transcendental. Here " $\alpha^{\beta "}$ " stands for $\exp (\beta \log \alpha)$, and any nonzero value of $\log \alpha$ is permissible. For a proof of the Gelfond-Schneider result, see, e.g., [Hua82, §17.9].

There are many sequences not discussed in $\S 10$ where it would be of interest to decide if they contain infinitely many primes, or composites. For example, fix a nonintegral rational number $\alpha>1$, and consider the sequence of numbers $\left\lfloor\alpha^{n}\right\rfloor$. Whiteman has conjectured that this sequence always contains infinitely many primes. If we drop the rationality condition, then from a very general theorem of Harman [Har97] we have that each sequence $\left\lfloor\alpha^{n}\right\rfloor$ contains infinitely many primes as long as $\alpha>1$ avoids a set of measure zero. (Of course since the rational numbers have measure zero, this has no direct consequence for Whiteman's conjecture.) Very little is known about the sequences considered by Whiteman. For the particular numbers $\alpha=3 / 2$ and $\alpha=4 / 3$, Forman \& Shapiro [FS67] present ingenious elementary arguments showing that the sequence $\left\lfloor\alpha^{n}\right\rfloor$ contains infinitely many composite numbers. Some extensions of their results have been obtained by Dubickas \& Novikas [DN05]; e.g., these authors prove that if $\xi>0$ and $\alpha \in\{2,3,4,6,3 / 2,4 / 3,5 / 4\}$, then the sequence $\left\lfloor\xi \alpha^{n}\right\rfloor$ contains infinitely many composites.

## Exercises

1. (Harris [Har56]) Let $b_{0}, b_{1}, b_{2}$ be positive integers with $b_{0}$ coprime to $b_{2}$. Define $A_{k}$ for $k=0,1$ and 2 as the numerator when the finite continued fraction

$$
b_{0}+\frac{1}{b_{1}+\frac{1}{\ddots+\frac{1}{b_{k}}}}
$$

is put in lowest terms. For $k=3,4, \ldots$, inductively define $b_{k}$ and $A_{k}$ by

$$
b_{k}=A_{0} A_{1} \cdots A_{k-3}
$$

and $A_{k}$ by the rule given above. Prove that the $A_{i}$ form an increasing sequence of pairwise coprime positive integers.
2. (Aldaz \& Bravo [AB03]) Let $p_{i}$ denote the $i$ th prime. Euclid's argument shows that for each $r$, there is a prime in the interval $\left(p_{r}, \prod_{1}^{r} p_{i}+1\right]$. Prove that the number of primes in the (smaller) interval ( $\left.p_{r}, \prod_{2}^{r} p_{i}+1\right]$ tends to infinity with $r$. Suggestion: With $P=\prod_{2}^{r} p_{i}$, show that $P$ $2, P-2^{2}, \ldots, P-2^{k}$ are $>1$ and pairwise coprime for fixed $k$ and large $r$; then choose a prime factor of each.
3. (Chowdhury [Cho89]) It is trivial that for $n \geq 1$, the number $n!+1$ has a prime divisor exceeding $n$. Show that for $n \geq 6$, the same holds for each of the numbers $n!+k$, where $2 \leq k \leq n$.
4. (Hegyvári [Heg93]) Suppose $a_{1}<a_{2}<a_{3}<\ldots$ is an increasing sequence of natural numbers for which $\sum 1 / a_{i}$ diverges. Show that the real number $\alpha:=0 . a_{1} a_{2} a_{3} \ldots$ formed by concatenating the decimal expansions of the $a_{i}$ is irrational. In particular, $0.235711131719 \ldots$ is irrational. Hint: First show that every finite sequence of decimal digits appears in the expansion of $\alpha$.

Remark. Suppose that in place of our divergence hypotheses, we assume that for each fixed $\theta<1$, the number of $a_{i} \leq x$ exceeds $x^{\theta}$ for all sufficiently large $x$. Then Copeland \& Erdős [CE46] have proved that the number $\alpha$ constructed above is normal (in base 10); in other words, not only does every finite digit string appear in the expansion of $\alpha$, but each string of length $k$ appears with the expected frequency $10^{-k}$.
5. (Euler) In courses in complex analysis, it is often proved that $\sin x$ possesses the following Weierstrass factorization (valid for all $x \in \mathbf{C}$ ):

$$
\begin{equation*}
\sin x=x \prod_{n=1}^{\infty}\left(1-\frac{x^{2}}{n^{2} \pi^{2}}\right) \tag{1.15}
\end{equation*}
$$

see, e.g., [Pri01] for a short, direct proof of this identity. A proof using only real-variable methods appears in [Kob84, Chapter II].
(a) Starting from (1.15), show that

$$
x \cot x=1-2 \sum_{m=1}^{\infty} \zeta(2 m) \frac{x^{2 m}}{\pi^{2 m}},
$$

where $\zeta$ denotes the Euler-Riemann zeta function. Hint: Take the logarithmic derivative of both sides.
(b) Computing by hand the first few coefficients in the Taylor series for $x \cot x$ about $x=0$, check that $\zeta(2)=\pi^{2} / 6$ and $\zeta(4)=\pi^{4} / 90$.
6. (J. D. Dixon) We outline Dixon's proof [Dix62] that $\pi$ is not the root of a polynomial over $\mathbf{Z}$ of degree $\leq 2$. The method is that employed by Niven to show $\pi$ is irrational (see $[\operatorname{Niv} 47]$ ). Suppose for the sake of contradiction that $\pi$ is a root of $P(T)=a T^{2}+b T+c$, where $a, b$ and $c$ are integers, not all vanishing.

Given a polynomial $f(T) \in \mathbf{R}[T]$, define

$$
\begin{equation*}
F(T):=f(T)-f^{(2)}(T)+f^{(4)}(T)-f^{(6)}(T)+\cdots \tag{1.16}
\end{equation*}
$$

Then $F(T) \in \mathbf{R}[T]$. View $F$ as a function of a real variable $x$.
(a) Check that

$$
\frac{d}{d x}\left(F^{\prime}(x) \sin x-F(x) \cos x\right)=f(x) \sin (x)
$$

and conclude that

$$
\begin{equation*}
\int_{0}^{\pi} f(x) \sin x d x=F(\pi)+F(0) \tag{1.17}
\end{equation*}
$$

(b) With $n$ a positive integer to be chosen shortly, let $f$ be the polynomial

$$
f(T):=\frac{1}{n!} P(T)^{2 n}(P(T)-P(0))^{2 n}
$$

Show that the left-hand side of (1.17) is strictly between 0 and 1 if $n$ is sufficiently large.
We now fix such an $n$ and derive a contradiction by showing that the right-hand side of (1.17) is an integer.
(c) Show that $f^{(r)}(0)=f^{(r)}(\pi)=0$ for all $0 \leq r<2 n$.
(d) If $e$ and $r$ are nonnegative integers and $r$ is even, show that there is an expansion of the form

$$
\frac{d^{r}}{d x^{r}}\left(P(x)^{e}\right)=\sum_{j=r / 2}^{r} c_{j} j!\binom{e}{j} P(x)^{e-j}
$$

for certain integers $c_{j}$.
(e) Use the result of part (d) to show that if $e$ is a nonnegative integer and $r \geq 2 n$ is even, then $\frac{1}{n!} \frac{d^{r}}{d x^{r}}\left(P(x)^{e}\right)$ is a polynomial in $P(x)$ with integer coefficients. Conclude that $f^{(r)}(0)$ and $f^{(r)}(\pi)$ are integers.
(f) Referring back to definition (1.16), deduce that $F(\pi)+F(0) \in \mathbf{Z}$.
7. In this exercise we present a proof similar to that of J. Hacks (on p. 8) but relying on the irrationality of $\pi$ in place of $\pi^{2}$. Let

$$
\chi(n)= \begin{cases}(-1)^{(n-1) / 2} & \text { if } 2 \nmid n \\ 0 & \text { otherwise }\end{cases}
$$

a) Show that $\chi(n)$ is a completely multiplicative function, i.e.,

$$
\chi(a b)=\chi(a) \chi(b)
$$

for every pair of positive integers $a, b$.
b) Assume that there are only finitely many primes. Show that for every $s>0$,

$$
\sum_{n=1}^{\infty} \frac{\chi(n)}{n^{s}}=\prod_{p}\left(1-\frac{\chi(p)}{p^{s}}\right)^{-1}
$$

c) Take $s=1$ and obtain a contradiction to the irrationality of $\pi$. You may assume that $\frac{\pi}{4}=1-\frac{1}{3}+\frac{1}{5}-\frac{1}{7}+\cdots$.
8. Say that a natural number $n$ is squarefull if $p^{2} \mid n$ whenever $p \mid n$, i.e., if every prime showing up in the factorization of $n$ occurs with multiplicity $>1$. Every perfect power is squarefull, but there are many other examples, such as $864=2^{5} \cdot 3^{3}$. Using Theorem 1.2 , show that $\sum^{\prime} n^{-1}$ converges to $\frac{\zeta(2) \zeta(3)}{\zeta(6)}$, where the ' indicates that the sum is restricted to squarefull $n$. Determine the set of real $\alpha$ for which $\sum^{\prime} n^{-\alpha}$ converges.
9. (Continuation) Show that every squarefull number has a unique representation in the form $u^{2} v^{3}$, where $u$ and $v$ are positive integers with $v$ squarefree. Deduce that for $x \geq 1$,

$$
\sum_{\substack{n \leq x \\ n \text { squarefull }}} 1=\frac{\zeta(3 / 2)}{\zeta(3)} x^{1 / 2}+O\left(x^{1 / 3}\right)
$$

10. (Ramanujan) Assuming $\zeta(2)=\pi^{2} / 6$ and $\zeta(4)=\pi^{4} / 90$, show that

$$
\sum^{\prime} \frac{1}{n^{2}}=\frac{9}{2 \pi^{2}},
$$

where the ' indicates that the sum ranges over positive squarefree integers $n$ with an odd number of prime divisors.
11. (Cf. Porubský [Por01]) If $R$ is a commutative ring, its Jacobson radical $J(R)$ is the intersection of all of its maximal ideals. Show that

$$
J(R)=\{x \in R: 1-x y \text { is invertible for all } y \in R\} .
$$

Deduce that if $R$ is an integral domain with finitely many units, then $J(R)=\{0\}$. Use this to prove that if $R$ is a principal ideal domain with finitely many units, then either $R$ is a field or $R$ contains an infinite set of pairwise nonassociated primes.
12. By carefully examining the proof of Theorem 1.10 , show that the theorem remains correct when $A=1$, provided that in condition (ii) we replace "prime" with "prime or equal to 1 ".
13. Suppose that $a_{1}<a_{2}<a_{3}<\cdots$ is an increasing sequence of natural numbers, and put $A(x):=\sum_{a_{i} \leq x} 1$. Prove that if $(\log x)^{-k} A(x) \rightarrow \infty$ for each fixed $k$, then infinitely many primes $p$ divide some $a_{i}$. Use this to give another proof of Lemma 1.17.
14. Prove the following theorem of Bauer [Bau06]:

Theorem. If $F(T) \in \mathbf{Z}[T]$ is a nonconstant polynomial with at least one real root, then for every $m \geq 3$, there exist infinitely many prime divisors $p$ of $F$ with $p \not \equiv 1(\bmod m)$.

Proceed by showing that each of the following conditions on $F$ is sufficient for the conclusion of the theorem to hold:
(a) $F$ has a positive leading coefficient and constant term -1 .
(b) $F$ has a positive leading coefficient and negative constant term.
(c) $F$ has a positive leading coefficient and $F(a)<0$ for some $a \in \mathbf{Z}$.
(d) $F$ has a positive leading coefficient and $F(a)<0$ for some $a \in \mathbf{Q}$.
(e) $F$ has a positive leading coefficient and $F(a)<0$ for some $a \in \mathbf{R}$.
(f) $F$ has a positive leading coefficient and $F(a)=0$ for some $a \in \mathbf{R}$. Hint for $(f)$ : Reduce to the case when $F$ has no multiple roots.
15. Let $F$ be a field of characteristic not dividing $m$. By carefully examining the proof of Lemma 1.19, show that the roots of $\Phi_{m}(T)$ in the algebraic closure of $F$ are precisely the primitive $m$ th roots of unity there, and that all these roots are simple.
16. (Continuation; Kronecker [Kro88], Dirichlet, Bauer [Bau06]) Define $\Phi_{m}(X, Y)$ as the homogenization of $\Phi_{m}(T)$, so that

$$
\Phi_{m}(X, Y)=\prod_{\substack{\zeta^{m}=1 \\ \zeta^{j} \neq 1 \text { if } 1 \leq j<m}}(X-\zeta Y)
$$

(a) Suppose $m>2$. Show that $\Phi_{m}(X+Y, X-Y)=G_{m}\left(X, Y^{2}\right)$ for some polynomial $G_{m}$ (say) with integer coefficients. Show also that $\prod_{d \mid m} d^{\mu(m / d)}$ is the coefficient of $X^{\varphi(m)}$ in $\Phi_{m}(X+Y, X-Y)$.
(b) Let $F$ be a field of characteristic not dividing $m$. Suppose $s$ is a nonsquare integer, and let $\sqrt{s}$ denote a fixed square root of $s$ from the algebraic closure of $F$. Show that the roots of $G_{m}(T, s) \in \mathbf{Z}[T]$ in the algebraic closure of $F$ are precisely the elements

$$
\sqrt{s} \frac{\zeta+1}{\zeta-1}
$$

where $\zeta$ runs through the primitive $m$ th roots of unity.
(c) Suppose $s$ is as in (b), and let $p$ be a prime for which $p \nmid 2 m s$. Show that $p$ is a prime divisor of $G_{m}(T, s)$ if and only if $p \equiv\left(\frac{s}{p}\right)(\bmod m)$.
(d) Show that if $p \equiv-1(\bmod 4)$ is a prime divisor of $G_{m}(T,-1)$ which does not divide $m$, then $p \equiv-1(\bmod m)$. Use Exercise 14 to show that $G_{m}(T,-1)$ has infinitely many such prime divisors, and deduce that there are infinitely many primes $p \equiv-1(\bmod m)$.
17. (M. Hirschhorn [Hir02]) Let $p_{1}<p_{2}<p_{3}<\cdots$ denote the sequence of odd primes.
(a) Let $N \in \mathbf{N}$. Prove that the number of odd positive integers $\leq N$ which can be written in the form $p_{1}^{e_{1}} \cdots p_{k}^{e_{k}}$ does not exceed

$$
\prod_{i=1}^{k}\left(\frac{\log N}{\log p_{i}}+1\right)<\left(\log \left(p_{k} N\right)\right)^{k}<\sqrt{2 k!} \sqrt{p_{k} N}
$$

Hint: Show that $(\log u)^{k} u^{-1 / 2} \leq(2 k / e)^{k}$ whenever $u \geq 1$. Now invoke the inequality $m!\geq(m / e)^{m}$, valid for every integer $m \geq 0$.
(b) Supposing that $p_{1}, \ldots, p_{k}$ exist (i.e., that there are at least $k$ odd primes), prove that $p_{k+1}$ exists and satisfies $p_{k+1} \leq 4(2 k!) p_{k}+1$.
18. Suppose that $A$ is a commutative monoid (written multiplicatively) and that $P$ is a system of generators for $A$, so that each element of $A$ can be written in the form $\prod_{p \in P} p^{e_{p}}$, where each $e_{p} \geq 0$ and only finitely many of the $e_{p}$ are nonzero. (We do not require that this representation be unique.) Suppose also that there is a function $\|\cdot\|: A \rightarrow \mathbf{N}$ with the following two properties:
(a) $\|\cdot\|$ respects multiplication, i.e., $\|a b\|=\|a\|\|b\|$ for all $a, b \in A$.
(b) For some real number $x_{0}$ and constants $c_{1}, c_{2}>0$, we have

$$
\begin{equation*}
c_{1} x \leq \#\{a \in A:\|a\| \leq x\} \leq c_{2} x \quad \text { for all } \quad x>x_{0} . \tag{1.18}
\end{equation*}
$$

Prove that $P$ is infinite, and that in fact $\sum_{p \in P} \frac{1}{\|p\|}$ diverges.
19. (Continuation)
(a) For each nonzero Gaussian integer $\alpha$ put $\|\alpha\|=|\alpha|^{2}$. Show that $\sum_{\pi}\|\pi\|^{-1}$ diverges, where the sum is over all Gaussian primes $\pi$. Deduce that $\sum_{p \equiv 1(\bmod 4)} p^{-1}$ diverges, where the sum is over rational primes $p \equiv 1 \bmod 4$.
(b) For each nonzero polynomial $F(T) \in \mathbf{F}_{q}[T]$, put $\|F\|:=q^{\operatorname{deg} F}$. Show that $\sum\|P\|^{-1}$ diverges, where $P$ ranges over the irreducible elements of $\mathbf{F}_{q}[T]$.
20. This exercise outlines a proof of Theorem 1.21 via algebraic number theory. Let $m$ be a positive integer, and let $\zeta$ be a primitive $m$ th root of unity. Put $K=\mathbf{Q}\left(\zeta_{m}\right)$, and identify $\operatorname{Gal}(K / \mathbf{Q})$ with $(\mathbf{Z} / m \mathbf{Z})^{\times}$. Let $H$ be a subgroup of $(\mathbf{Z} / m \mathbf{Z})^{\times}$, and let $L \subset K$ be the fixed field of $H$.
(a) Say that two sets of rational primes $\mathcal{P}_{1}$ and $\mathcal{P}_{2}$ eventually coincide if their symmetric difference is finite; in this case we write $\mathcal{P}_{1} \doteq$ $\mathcal{P}_{2}$. Prove that $\mathcal{P}_{1} \doteq \mathcal{P}_{2}$, where $\mathcal{P}_{1}$ is the set of primes for which $p \bmod m \in H$ and $\mathcal{P}_{2}$ is the set of primes which split completely in L. Hint: If $p$ is a prime not dividing $m$, analyze how the Frobenius element of $p$ in $\operatorname{Gal}(K / \mathbf{Q})$ behaves upon restriction to $L$.
(b) Let $\theta$ be an algebraic integer for which $L=\mathbf{Q}(\theta)$. Let $F$ be the minimal polynomial of $\theta$. Prove that $\mathcal{P}_{2}$, and hence also $\mathcal{P}_{1}$, eventually coincides with the set of prime divisors of $F$. Hint: $L / \mathbf{Q}$ is Galois, so an unramified rational prime splits completely in $L$ exactly when it has a degree 1 prime factor; now apply the Kummer-Dedekind theorem.
21. (Pólya [Pól21]; see also [MS00]) Suppose that $a$ and $b$ are nonzero integers and $a \neq \pm 1$. Let $\mathcal{P}$ be the set of primes for which the exponential congruence $a^{k} \equiv b(\bmod p)$ has a positive integer solution $k$. In other words, $\mathcal{P}$ is the set of primes which divide some term of the sequence

$$
a-b, \quad a^{2}-b, \quad a^{3}-b, \quad a^{4}-b, \ldots
$$

This exercise outlines a proof that $\mathcal{P}$ is always an infinite set.
We may suppose that $b$ is not a power of $a$, as otherwise $\mathcal{P}$ contains every prime. We assume for the sake of contradiction that $\mathcal{P}$ is finite.
(a) For each $p \in \mathcal{P}$ and each $k \geq 1$, define integers $v_{p, k} \geq 0$ by writing

$$
a^{k}-b= \pm \prod_{p \in P} p^{v_{p, k}}
$$

For each $p \in \mathcal{P}$, set $v_{p}:=\sup _{k \geq 1} v_{p, k}$. We let $\mathcal{P}_{1}:=\left\{p \in \mathcal{P}: v_{p}<\right.$ $\infty\}$ and we put $\mathcal{P}_{2}:=\mathcal{P} \backslash \mathcal{P}_{1}$. Show that if $p \in \mathcal{P}_{2}$, then $p \nmid a$.
(b) Suppose $p \in \mathcal{P}_{2}$, and let $l_{p}$ be the order of $a$ modulo $p$. (This exists by part (a).) Define $e_{p}$ so that $p^{e_{p}} \| a^{l_{p}}-1$. Show that if $k$ is a positive integer for which $p^{e_{p}+1} \mid a^{k}-b$, then $k$ belongs to a fixed residue class modulo $p$.
(c) Show that there is an infinite arithmetic progression of integers $k$ which avoid all the residue classes $\bmod p\left(p \in \mathcal{P}_{2}\right)$ determined in (b). Prove that $a^{k}-b$ is uniformly bounded for such $k$, contradicting that $\left|a^{k}-b\right| \rightarrow \infty$ as $k \rightarrow \infty$.

Remark. In the opposite direction, one can ask when the set $\mathcal{P}$ defined above omits infinitely many primes. Using the Chebotarev density theorem, Schinzel $[\mathbf{S c h 6 0}]$ has shown that this holds unless $b=a^{k}$ for some nonnegative integer $k$. See also [MS00].
22. (Křížek et al. $\left[\mathbf{K L S 0 2 ]}\right.$ ) Let $F_{n}=2^{2^{n}}+1$ be the $n$th Fermat number. Suppose $N \in \mathbf{N}$.
(a) Show that there are fewer than $2^{N}$ distinct prime divisors of the product $F_{0} \cdots F_{N-1}$.
(b) Show that for each $x>0$, the number of primes $p \leq x$ which divide $F_{n}$ for some $n \geq N$ is at most $x / 2^{N+1}$.
(c) Making an appropriate choice of $N$, deduce from (a) and (b) that there are $\ll \sqrt{x}$ primes $p \leq x$ which divide a term of the sequence $F_{0}, F_{1}, F_{2}, \ldots$
(d) Deduce that if $\lambda>1 / 2$, then $\sum^{\prime} p^{-\lambda}<\infty$, where the ${ }^{\prime}$ indicates that the sum is restricted to primes dividing at least one Fermat number. When $\lambda=1$, this confirms a conjecture of Golomb [Gol55].
23. (Erdős \& Turán [ET34]) For $n>1$, write $P(n)$ for the largest prime factor of $n$. In this exercise we show that if $S$ is an infinite set of natural numbers, then

$$
\begin{equation*}
\{P(a+b): a, b \in S\} \quad \text { is unbounded. } \tag{1.19}
\end{equation*}
$$

For each prime $p$, let $v_{p}$ be the $p$-adic valuation, defined so that $p^{v_{p}(n)} \| n$ for every natural number $n$.
(a) Let $S$ be an arbitrary infinite set of natural numbers. Show that for each odd prime $p$, we can determine an infinite subset $S^{\prime} \subset S$ with the property that whenever $a, b \in S^{\prime}$,

$$
v_{p}(a+b)=\min \left\{v_{p}(a), v_{p}(b)\right\} .
$$

Hint: First treat the case when no element of $S$ is divisible by $p$.
(b) Suppose, for the sake of contradiction, that $S$ is infinite but (1.19) fails. Using part (a), argue that we may assume (1.20) holds for

$$
\begin{aligned}
& 5^{n} 7^{d} 13 \\
& \downarrow(A B)^{d} J \\
& 2^{d} 3^{d} 5^{n-d} 11 \\
& \downarrow(E F)^{d} K \\
& 2^{d} 5^{n-d} 7^{d} 13 \\
& \downarrow(A B)^{d} J \\
& 2^{2 d} 3^{d} 5^{n-2 d} 11 \\
& \downarrow(E F)^{d} K \\
& 2^{2 d} 5^{n-2 d} 7^{d} 13 \\
& \downarrow(A B)^{d} J \\
& \downarrow(E F)^{d} K \\
& 2^{q d} 5^{r} 7^{d} 13 \\
& \downarrow(A B)^{r} A \\
& \text { (if } r>0) \stackrel{2^{n} 3^{r} 7^{d-r-1} 17}{\swarrow} \text { (if } r=0 \text { ) } \\
& 2^{n} 3^{r-1} 7^{d-r-1} 19 \quad 2^{n} 7^{d-1} \\
& \downarrow(D G)^{n} H \quad \downarrow L^{n} M^{d-1} N \\
& 3^{r-1} 5^{n} 7^{d-r} 113^{n-1} K 5^{n+1} 11 \\
& \downarrow(E F)^{r-1} K \quad \downarrow(E F)^{n} K \\
& 5^{n} 7^{d-1} 13 \quad 5^{n+1} 7^{n} 13
\end{aligned}
$$

Figure 4. The action of Conway's prime-producing machine when started with $5^{n} 7^{d} 13$, where $0<d<n$. The variables $q$ and $d$ are defined by the division algorithm: $n=d q+r$ where $0 \leq r<d$.
every pair $a, b \in S$ and every odd prime $p$. We make this assumption from now on.
(c) Now argue that $v_{2}(a)=v_{2}(b)$ for every pair of elements $a, b \in S$. Thus, dividing through by a suitable power of 2 , we may (and do) assume that all the elements of $S$ are odd.
(d) Finally, show that for each pair of elements $a, b \in S$, we have

$$
a+b=2^{v_{2}(a+b)} \prod_{p>2} p^{\min \left\{v_{p}(a), v_{p}(b)\right\}}
$$

Show that this equation leads to a contradiction if $a$ and $b$ are chosen to be congruent modulo 4.
24. Figure 4, based on Conway's article [Con87], describes the action of Conway's prime-producing machine. Decipher this figure and explain how it proves Theorem 1.8. For a more detailed explanation of the workings of Conway's prime-producing machine, see Guy's expository article [Guy83].
25. (Schinzel [Sch62a]) In 1857, Bunyakovsky conjectured [Bun57] that if $F(T) \in \mathbf{Z}[T]$ is an irreducible polynomial with positive leading coefficient and $D$ is the largest positive integer dividing $F(n)$ for each $n \in \mathbf{Z}$, then $F(n) / D$ is prime for infinitely many natural numbers $n$. Show that this would follow from Hypothesis H.
26. (Granville; see, e.g., [Mol97, Theorem 2.1]) Assume Hypothesis H. Show that for every natural number $N_{0}$, one can find a positive integer $A$ with the property that $n^{2}+n+A$ assumes prime values for all $0 \leq n \leq N_{0}$. Hint: Apply Hypothesis H to the $N_{0}$ linear polynomials $T, T+\left(1^{2}+1\right), T+\left(2^{2}+2\right), \ldots, T+\left(N_{0}^{2}+N_{0}\right)$.
27. (Schinzel \& Sierpiński [SS58]) Assume Hypothesis H. Show that if $n>1$ and $r$ is a positive integer divisible by all primes $p \leq n$, then there are infinitely many arithmetic progressions of length $n$ and common difference $r$ consisting of consecutive primes.
Remark. The weaker claim that there are arbitrarily long arithmetic progressions of primes was recently proved in a technical tour de force by Green \& Tao [GT08], using ideas borrowed from ergodic theory (and several other fields). For some striking elementary consequences of the Green-Tao result, see [Gra08a].
28. (Cf. Chang \& Lih [CL77]) Show that for every $N \in \mathbf{N}$, there is a polynomial $F(T) \in \mathbf{Z}[T]$ for which $\{F(k)\}_{k=0}^{N}$ is a sequence of $N+1$ distinct primes. Hint: For $0 \leq k \leq N$, put $c_{k}(T)=\prod_{0 \leq i \leq N, i \neq k}(T-$ $i)$. Using Corollary 1.20 , choose integers $r_{0}, r_{1}, \ldots, r_{N}$ for which $\{1+$ $\left.r_{k} c_{k}(k)\right\}_{k=0}^{N}$ is a sequence of $N+1$ distinct primes. Put $F(T):=1+$ $\sum_{i=0}^{N} r_{i} c_{i}(T)$.
29. (Clement [Cle49], Cucurezeanu [Cuc68]) Let $k$ and $n$ be integers with $n>k \geq 2$. Suppose that $n$ has no prime divisors $<k$. Show that $n$ and $n+k$ are simultaneously prime if and only if

$$
k \cdot k!((n-1)!+1)+\left(k!-(-1)^{k}\right) n \equiv 0(\bmod n(n+k)) .
$$

30. (Shanks $[$ Sha64 $])$ Let $F(z)=\sum_{n=0}^{\infty} z^{n(n+1) / 2}$ and define

$$
G(z):=(F(z)-1)^{2}-(F(z)-1) .
$$

Prove that there are infinitely many primes of the form $\frac{n^{2}+1}{2}$ (with $n \in \mathbf{N}$ ) if and only if the power series expansion of $G$ has infinitely many negative coefficients.
31. Suppose $p \equiv 3(\bmod 4)$ is prime. Prove that if $2 p+1$ is also prime, then $2 p+1 \mid 2^{p}-1$. Deduce that Hypothesis H implies Conjecture 1.27.
32. (Selfridge; cf. $[\mathbf{E r d 5 0 b}]$ ) Let $n \in \mathbf{N}$. Show that $78557 \cdot 2^{n}+1$ is divisible by some prime number from the set $\{3,5,7,13,19,37,73\}$. In particular, $78557 \cdot 2^{n}+1$ is always composite.

Table 1. Mann-Shanks criterion: Columns containing only bold entries are indexed by prime numbers.

|  | 0 | 1 | $\mathbf{2}$ | $\mathbf{3}$ | 4 | $\mathbf{5}$ | 6 | $\mathbf{7}$ | 8 | 9 | 10 | $\mathbf{1 1}$ | 12 | $\mathbf{1 3}$ |
| ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: |
| 0 | 1 |  |  |  |  |  |  |  |  |  |  |  |  |  |
| 1 |  |  | $\mathbf{1}$ | $\mathbf{1}$ |  |  |  |  |  |  |  |  |  |  |
| 2 |  |  |  |  | 1 | $\mathbf{2}$ | 1 |  |  |  |  |  |  |  |
| 3 |  |  |  |  |  |  | 1 | $\mathbf{3}$ | $\mathbf{3}$ | 1 |  |  |  |  |
| 4 |  |  |  |  |  |  |  |  | 1 | $\mathbf{4}$ | 6 | $\mathbf{4}$ | 1 |  |
| 5 |  |  |  |  |  |  |  |  |  |  | 1 | $\mathbf{5}$ | $\mathbf{1 0}$ | $\mathbf{1 0}$ |
| 6 |  |  |  |  |  |  |  |  |  |  |  | 1 | $\mathbf{6}$ |  |

33. (Louisiana State University Problem Solving Group [PSG02]) Prove that $5^{4 n}+5^{3 n}+5^{2 n}+5^{n}+1$ is composite for every natural number $n$. If you know some algebraic number theory, establish the following generalization: If $q>1$ is a squarefree natural number with $q \equiv$ $1(\bmod 4)$, then $\Phi_{q}\left(q^{n}\right)$ is composite for every natural number $n$.

Hint (due to J. A. Rouse): $q^{n}-\zeta$ is a difference of squares in $\mathbf{Z}[\zeta]$, where $\zeta$ denotes a primitive $q$ th root of unity.
34. Table 1 illustrates a primality criterion discovered by Mann \& Shanks [MS72]: Place the rows of Pascal's triangle in an infinite table, where the zeroth row (consisting of the single element 1) is placed in column 0 . Each successive row is shifted two units right. An element of the $n$th row is written in boldface when it is divisible by $n$. Then the column number is prime exactly when all entries in its column are written in boldface. Prove this!
35. (Hayes [Hay65]) Suppose that $R$ is a principal ideal domain with infinitely many prime ideals. Show that every nonconstant polynomial $A$ over $R$ can be written as the sum of two irreducible polynomials of the same degree as $A$. Hint: Arrange for both summands to satisfy the Eisenstein criterion with respect to the same prime.

## Cyclotomy

The principles upon which the division of the circle depend, and geometrical divisibility of the same into seventeen parts, etc. - C. F. Gauss

## 1. Introduction

The terse quotation opening this chapter also opens Gauss's mathematical diary, commenced on March 30, 1796, when Gauss was 18 years old. This entry carries more significance for mathematics than a straight reading would suggest; it was his discovery of the constructibility of the regular 17-gon that swayed Gauss to choose mathematics over philology, his other early love.

It has been known since the time of Euclid that the regular $n$-gon is constructible for any $n \geq 3$ of the form

$$
n=2^{a} 3^{b} 5^{c} \quad \text { where } \quad a \geq 0, \quad b=0 \text { or } 1, \quad c=0 \text { or } 1
$$

Whether there were other constructible regular polygons remained an open question for 2000 years. The millenia-long silence was broken by the following notice, which appeared in the April 1796 Allgemeine Literaturzeitung (see [Dun04, p. 28]):

It is known to every beginner in geometry that various regular polygons, viz., the triangle, tetragon, pentagon, 15-gon and those which arise by the continued doubling of the number of sides of one of them, are geometrically constructible.

One was already that far in the time of Euclid, and, it seems, it has generally been said since then that the field of
elementary geometry extends no farther: at least I know of no successful attempt to extend its limits on this side.

So much the more, methinks, does the discovery deserve attention. . that besides those regular polygons a number of others, e.g., the 17-gon, allow of a geometrical construction. This discovery is really only a special supplement to a theory of greater inclusiveness, not yet completed, and is to be presented to the public as soon as it has reached its completion.

Carl Friedrich Gauss
Student of Mathematics at Göttingen
This "theory of greater inclusiveness" (which became known as cyclotomy; literally, "circle-splitting") appeared five years later in the last of the seven sections of the Disquisitiones. There Gauss [Gau86, §365] offers a complete characterization of the constructible regular polygons. Recall that a Fermat prime is a prime number of the form $2^{n}+1$, where $n$ is a positive integer.

Theorem 2.1 (Gauss, Wantzel). It is possible to construct a regular n-sided polygon in the plane by straightedge and compass if and only if $n=2^{e} p_{1} \cdots p_{k}$ for $e \geq 0$ and distinct Fermat primes $p_{1}, \ldots, p_{k}$ (where $k \geq 0$ ).

Wantzel's name is attached to this result because the Disquisitiones, while insisting on the necessity of the condition of Theorem 2.1, proves only its sufficiency. The first published proof that the regular $n$-gon is constructible only for those $n$ as in Theorem 2.1 is due to Wantzel [Wan37].

The first goal of this chapter is to prove the Gauss-Wantzel theorem. The remainder of this chapter discusses two applications of cyclotomy to the study of reciprocity laws.

Recall that when $p$ is an odd prime and $a$ is an integer relatively prime to $p$, the Legendre symbol $\left(\frac{a}{p}\right)$ is defined to be 1 if $a$ is a square modulo $p$ and -1 otherwise. Gauss was the first to prove the following fundamental result, which to this day forms the capstone of many a course in elementary number theory:

Theorem 2.2 (Law of quadratic reciprocity). Suppose that $p$ and $q$ are distinct odd primes. Then

$$
\left(\frac{q}{p}\right)=(-1)^{\frac{p-1}{2} \frac{q-1}{2}}\left(\frac{p}{q}\right) .
$$

Over the course of his life Gauss worked out eight different proofs of Theorem 2.2. Eight proofs may seem like overkill, but Gauss was hoping that these arguments would shed light on the theory of higher power residues (cubic residues, quartic residues, etc.). Arguably the first significant step
in this direction came in September 1796, when Gauss found two proofs of Theorem 2.2, both based on cyclotomy.

In $\S 6$ we present a "cyclotomic" proof of quadratic reciprocity. To illustrate the scope of these methods, we turn next to a study of cubic residues. Notice that if $p$ is prime and $p \equiv 2(\bmod 3)$, then 3 is coprime to $p-1=\# \mathbf{F}_{p}^{\times}$, and so every element of $\mathbf{F}_{p}^{\times}$is a cube. So cubic residues are only interesting for primes $p \equiv 1(\bmod 3)$. To state our results for these primes, we need the following elementary lemma:

Lemma 2.3. Let $p \equiv 1(\bmod 3)$ be prime. Then there are integers $L$ and $M$, uniquely determined up to sign, for which $4 p=L^{2}+27 M^{2}$.

Proof. We first show that $p$ can be written in the form $a^{2}+a b+b^{2}$. Since $p \equiv 1(\bmod 3)$ and $(\mathbf{Z} / p \mathbf{Z})^{\times}$is a cyclic group, there is an element of order 3 in $(\mathbf{Z} / p \mathbf{Z})^{\times}$and hence an integer $r$ satisfying $r^{2}+r+1 \equiv 0(\bmod p)$. Let $x$ and $y$ run over pairs of integers with $0 \leq x \leq \sqrt{p}$ and $0 \leq y \leq \sqrt{p}$, and consider the difference $x-r y$ modulo $p$. There are $(\lfloor\sqrt{p}\rfloor+1)^{2}>p$ such pairs, and so by the Pigeonhole principle, we have $x_{1}-r y_{1} \equiv x_{2}-r y_{2}(\bmod p)$ for some $x_{1}, y_{1}, x_{2}, y_{2}$ with $\left(x_{1}, y_{1}\right) \neq\left(x_{2}, y_{2}\right)$ and $0 \leq x_{i}, y_{i}<\sqrt{p}$. Then with $a=x_{1}-x_{2}$ and $b=y_{1}-y_{2}$, we have $(a, b) \neq(0,0), a \equiv r b(\bmod p)$ and $|a|,|b|<\sqrt{p}$. Moreover,

$$
a^{2}+a b+b^{2} \equiv\left(r^{2}+r+1\right) b^{2} \equiv 0(\bmod p) \quad \text { and } \quad 0<\left|a^{2}+a b+b^{2}\right|<3 p .
$$

So $a^{2}+a b+b^{2}=p$ or $a^{2}+a b+b^{2}=2 p$. Working modulo 2 , we see that $a^{2}+a b+b^{2}$ is even only when both $a$ and $b$ are even, in which case $a^{2}+a b+b^{2}$ is a multiple of 4 . Since $4 \nmid 2 p$, we must have $a^{2}+a b+b^{2}=p$, as desired.

If $b$ is a multiple of 3 , say $b=3 M$, then the lemma follows quickly: From $p=a^{2}+a b+b^{2}$ we deduce $4 p=(2 a+b)^{2}+3 b^{2}=(2 a+b)^{2}+27 M^{2}$. So we have the lemma with this value of $M$ and $L:=2 a+b$. By the symmetry in $a$ and $b$, the lemma also holds if $a$ is a multiple of 3 . So we can suppose that $3 \nmid a b$. In this case, from $a^{2}+a b+b^{2} \equiv p \equiv 1(\bmod 3)$ we deduce that $a b \equiv-1(\bmod 3)$, which forces $a \equiv-b(\bmod 3)$. Put $A=-b$ and $B=a+b$. Then $A^{2}+A B+B^{2}=a^{2}+a b+b^{2}=p$; moreover, 3 divides $B$, and so we can run our previous argument.

We leave the proof of uniqueness as Exercise 1.

It turns out that the numbers $L$ and $M$ play a pivotal role in the study of cubic residues modulo $p$. This is already evident in Table 1; a bit of staring at this table prompts the following guess:

Table 1. The first fifty primes $p \equiv 1(\bmod 3)$ together with positive values of $L$ and $M$ for which $4 p=L^{2}+27 M^{2}$ and the cubic residue status of 2 and 3 .

| $p$ | $L$ | $M$ | $2=$ cube? | $3 ?$ | $p$ | $L$ | $M$ | $2=$ cube? | $3 ?$ |
| ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: |
| 7 | 1 | 1 | N | N | 271 | 29 | 3 | N | Y |
| 13 | 5 | 1 | N | N | 277 | 26 | 4 | Y | N |
| 19 | 7 | 1 | N | N | 283 | 32 | 2 | Y | N |
| 31 | 4 | 2 | Y | N | 307 | 16 | 6 | Y | Y |
| 37 | 11 | 1 | N | N | 313 | 35 | 1 | N | N |
| 43 | 8 | 2 | Y | N | 331 | 1 | 7 | N | N |
| 61 | 1 | 3 | N | Y | 337 | 5 | 7 | N | N |
| 67 | 5 | 3 | N | Y | 349 | 37 | 1 | N | N |
| 73 | 7 | 3 | N | Y | 367 | 35 | 3 | N | Y |
| 79 | 17 | 1 | N | N | 373 | 13 | 7 | N | N |
| 97 | 19 | 1 | N | N | 379 | 29 | 5 | N | N |
| 103 | 13 | 3 | N | Y | 397 | 34 | 4 | Y | N |
| 109 | 2 | 4 | Y | N | 409 | 31 | 5 | N | N |
| 127 | 20 | 2 | Y | N | 421 | 19 | 7 | N | N |
| 139 | 23 | 1 | N | N | 433 | 2 | 8 | Y | N |
| 151 | 19 | 3 | N | Y | 439 | 28 | 6 | Y | Y |
| 157 | 14 | 4 | Y | N | 457 | 10 | 8 | Y | N |
| 163 | 25 | 1 | N | N | 463 | 23 | 7 | N | N |
| 181 | 7 | 5 | N | N | 487 | 25 | 7 | N |  |
| 193 | 23 | 3 | N | Y | 499 | 32 | 6 | N | N |
| 199 | 11 | 5 | N | N | 523 | 43 | 3 | Y | Y |
| 211 | 13 | 5 | N | N | 541 | 29 | 7 | N | Y |
| 223 | 28 | 2 | Y | N | 547 | 1 | 9 | N | N |
| 229 | 22 | 4 | Y | N | 571 | 31 | 7 | N | Y |
| 241 | 17 | 5 | N | N | 577 | 11 | 9 | N | N |

Theorem 2.4 (Gauss [Gau73a, §4]). Let $p \equiv 1(\bmod 3)$, and write $4 p=$ $L^{2}+27 M^{2}$, where $L$ and $M$ are positive. Then

2 is a cube $\bmod p \Longleftrightarrow 2 \mid L$ and $2 \mid M$

$$
\Longleftrightarrow p=L^{\prime 2}+27 M^{\prime 2} \text { for some } L^{\prime}, M^{\prime},
$$

and
3 is a cube $\bmod p \Longleftrightarrow 3 \mid M \Longleftrightarrow 4 p=L^{\prime 2}+243 M^{\prime 2}$ for some $L^{\prime}, M^{\prime}$.
We have labeled this in the style of a theorem, and indeed our guess can be proved correct. We will do this in $\S 7.2$.

Table 2. Primes $p \equiv 1(\bmod 3)$ between $10^{6}$ and $10^{6}+10^{3}$, together with the cubic residue status of $p$ with respect to 5,7 and 11 , and the ratios $\frac{L}{3 M}$ with respect to the same moduli.

| $p$ | $L$ | $M$ | $5 ?$ | $\frac{L}{3 M} \bmod 5$ | $7 ?$ | $\frac{L}{3 M} \bmod 7$ | 11? | $\frac{L}{3 M} \bmod 11$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 100003 | 337 | 103 | N | -2 | N | 1 | N | -4 |
| 100057 | 175 | 117 | Y | 0 | Y | 0 | N | 1 |
| 100069 | 458 | 84 | N | -1 | Y | $\infty$ | N | 4 |
| 100129 | 562 | 56 | N | -1 | Y | $\infty$ | N | 4 |
| 100153 | 443 | 87 | N | -2 | N | 1 | N | -1 |
| 100183 | 383 | 97 | N | -2 | N | 3 | N | 4 |
| 100189 | 209 | 115 | Y | $\infty$ | N | 3 | Y | 0 |
| 100207 | 421 | 91 | N | 2 | Y | $\infty$ | N | 4 |
| 100213 | 575 | 51 | Y | 0 | N | -1 | N | -3 |
| 100237 | 194 | 116 | N | -2 | N | 1 | N | 1 |
| 100267 | 224 | 114 | N | 2 | Y | 0 | N | 4 |
| 100279 | 137 | 119 | N | 1 | Y | $\infty$ | N | 1 |
| 100291 | 491 | 77 | N | 1 | Y | $\infty$ | Y | $\infty$ |
| 100297 | 250 | 112 | Y | 0 | Y | $\infty$ | Y | 5 |
| 100333 | 515 | 71 | Y | 0 | N | -1 | Y | 5 |
| 100357 | 631 | 11 | N | 2 | N | 3 | Y | $\infty$ |
| 100363 | 355 | 101 | Y | 0 | N | -1 | Y | -5 |
| 100393 | 593 | 43 | N | 2 | N | -3 | N | 4 |
| 100411 | 179 | 117 | N | -1 | N | -3 | N | -3 |
| 100417 | 139 | 119 | N | 2 | Y | $\infty$ | N | -3 |
| 100447 | 404 | 94 | N | 2 | N | -1 | N | -2 |
| 100459 | 263 | 111 | N | 1 | N | 1 | N | -4 |
| 100483 | 8 | 122 | N | -2 | N | -3 | N | -1 |
| 100501 | 323 | 105 | Y | $\infty$ | Y | $\infty$ | N | -1 |
| 100519 | 523 | 69 | N | -1 | N | 3 | N | -3 |
| 100537 | 305 | 107 | Y | 0 | N | 3 | N | 4 |
| 100549 | 83 | 121 | N | 1 | N | 1 | Y | $\infty$ |
| 100591 | 181 | 117 | N | 1 | N | -1 | Y | -5 |
| 100609 | 622 | 24 | N | 1 | N | 3 | N | 1 |
| 100621 | 574 | 52 | N | -1 | Y | 0 | N | 1 |
| 100669 | 626 | 20 | Y | $\infty$ | N | -1 | N | 2 |
| 100693 | 475 | 81 | Y | 0 | N | -3 | N | 2 |
| 100699 | 143 | 119 | N | -1 | Y | $\infty$ | Y | 0 |
| 100741 | 509 | 73 | N | 1 | N | -1 | N | -3 |
| 100747 | 605 | 37 | Y | 0 | N | -3 | Y | 0 |
| 100801 | 254 | 112 | N | -1 | Y | $\infty$ | N | 2 |
| 100927 | 380 | 98 | Y | 0 | Y | $\infty$ | N | -2 |
| 100957 | 185 | 117 | Y | 0 | N | 3 | N | 2 |
| 100981 | 457 | 85 | Y | $\infty$ | N | 3 | N | 3 |
| 100987 | 595 | 43 | Y | 0 | Y | 0 | N | -4 |
| 100999 | 452 | 86 | N | -1 | N | 3 | N | -2 |

Encouraged by this success, let us attempt to characterize the primes $p$ for which $q=5,7$ and 11 are cubic residues. Table 2 shows the results of a computation for primes $p \equiv 1(\bmod 3)$ between $10^{6}$ and $10^{6}+10^{3}$. This range of primes was motivated by the desire to see reasonably large values of $L$ and $M$. In this table we also include the ratio $\frac{L}{3 M} \bmod q$, writing $\infty$ for $\frac{L}{3 M} \bmod q$ when $q \mid M$. (Granted, it requires prophetic insight even to consider the ratio of $L$ to $M \bmod q$, and a double portion of such to consider the more obscure $\frac{L}{3 M}$. Patience; all will be clear in time!)

For $q=3,5$ and 7 , it appears from Table 2 that $q$ is a cube modulo $p$ precisely when $q \mid L M$ (i.e., when $\frac{L}{3 M}=0$ or $\infty$ ). When $q=11$, it seems that $q$ is a cube modulo $p$ if $\frac{L}{3 M}=0$ or $\infty$, but also when $\frac{L}{3 M}= \pm 5$. These limited examples lead us to conjecture that a fixed prime $q$ is a cubic residue of $p$ if and only if $\frac{L}{3 M} \bmod q$ belongs to a certain subset $S$ of $\mathbf{Z} / q \mathbf{Z} \cup\{\infty\}$.

We now state Jacobi's cubic reciprocity law, which vindicates our conjecture and provides an explicit description of the set $S$ :

Theorem 2.5 (Jacobi's cubic reciprocity law). Let $p$ and $q$ be distinct primes greater than 3 , and suppose that $p \equiv 1(\bmod 3)$. Jacobi:

$$
q \text { is a cube modulo } p \Longleftrightarrow \frac{L+3 M \sqrt{-3}}{L-3 M \sqrt{-3}} \text { is a cube in } \mathbf{F}_{q}(\sqrt{-3}) \text {. }
$$

Z.-H. Sun: Equivalently (as shown in detail in §7.4), let $G=G(q)$ be the group

$$
\left\{[a, b]: a, b \in \mathbf{F}_{q} \text { and } a^{2}+3 b^{2} \neq 0\right\}
$$

where $[a, b]$ and $[c, d]$ are identified if one is a nonzero scalar multiple of the other, and where multiplication is defined by

$$
[a, b] \odot[c, d]=[a c-3 b d, a d+b c] .
$$

Then $G$ is a cyclic group of order $q-\left(\frac{-3}{q}\right)$, and

$$
q \text { is a cube modulo } p \Longleftrightarrow[L, 3 M] \text { is a cube in } G \text {. }
$$

One can use Theorem 2.5 to compute $S$ for any given prime $q$. For the primes $q \leq 37$, this was carried out by Jacobi ([Jac27]; cf. [Jac69]); his results for $q=11,13,17,23,29,31$ and 37 are quoted in Table 3. (Jacobi considers the expression $\frac{L}{M}$ instead of $\frac{L}{3 M}$, but as we shall see in the proof, the latter arises somewhat more naturally.) We note that Jacobi's law appears (without proof) in Gauss's Nachlass [Gau73a, §2].

## 2. An algebraic criterion for constructibility

Let us review the rudiments of straightedge and compass constructions. (We assume a prior casual acquaintance with these of the type formed in a typical secondary-school geometry course; alternatively, all we need and more can

Table 3. Jacobi's criteria for $q=11,13,17,23,29,31$ or 37 to be cubic residues modulo $p=\frac{1}{4}\left(L^{2}+27 M^{2}\right)$. In each case it is necessary and sufficient that either $q|L, q| M$, or that one of the given congruences holds.

| $q$ | 11 | 13 | 17 | 19 | 23 | 29 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | $L \equiv \pm 4 M$ | $L \equiv \pm M$ | $L \equiv \pm 3 M$ | $L \equiv \pm 3 M$ | $L \equiv \pm 2 M$ | $L \equiv \pm 2 M$ |
|  |  |  | $L \equiv \pm 9 M$ | $L \equiv \pm 9 M$ | $L \equiv \pm 8 M$ | $L \equiv \pm M$ |
|  |  |  |  |  | $L \equiv \pm 11 M$ | $L \equiv \pm 11 M$ |
|  |  |  |  |  |  | $L \equiv \pm 13 M$ |
|  |  |  |  |  | 31 | 37 |
|  |  |  |  |  | $L \equiv \pm 5 M$ | $L \equiv \pm 8 M$ |
|  |  |  |  |  | $L \equiv \pm 7 M$ | $L \equiv \pm 3 M$ |
|  |  |  |  |  | $L \equiv \pm 6 M$ | $L \equiv \pm 9 M$ |
|  |  |  |  |  | $L \equiv \pm 11 M$ | $L \equiv \pm 7 M$ |
|  |  |  |  |  |  | $L \equiv \pm 12 M$ |

be found in the book of Courant \& Robbins [CR41, Chapter III, Part I].) We begin with two "constructed points" $O=(0,0)$ and $P=(0,1)$ in the plane $\mathbf{R}^{2}$. There are now three fundamental constructions we can perform:
(i) Given two constructed points, draw the line between them.
(ii) Given two constructed points, draw the line segment between them.
(iii) Given a constructed point and a constructed line segment, draw the circle centered at the given point with radius the length of the specified segment.

Each time two distinct lines intersect, or a line and a circle intersect, we add the point(s) of intersection to our set of constructible points. These processes may be continued indefinitely.

The key to proving Theorem 2.1 is to translate "constructibility" into an algebraic notion. Call $x+i y \in \mathbf{C}$ constructible if the point $(x, y) \in \mathbf{R}^{2}$ is constructible (in finitely many steps). Then one can prove:
Lemma 2.6. The complex number $\alpha$ is constructible if and only if there is a tower of subfields of the complex numbers

$$
\mathbf{Q}:=K_{0} \subset K_{1} \subset \cdots \subset K_{m},
$$

where $\alpha \in K_{m}$ and, for each $1 \leq i \leq m, K_{i}=K_{i-1}\left(\sqrt{\beta_{i}}\right)$ for some $\beta_{i} \in$ $K_{i-1}$. The set of constructible complex numbers forms a field under complex addition and multiplication.

We leave the proof of Lemma 2.6 as Exercise 4.
Lemma 2.6 reduces the Gauss-Wantzel theorem (Theorem 2.1) to an assertion in field theory and allows us to quickly dispense with the necessity
half of this result. We take for granted the (easy) fact that the constructibility of the $n$-gon is equivalent to the constructibility of an arbitrary primitive $n$th root of unity $\zeta_{n}$ (Exercise 5) and the fact that the cyclotomic polynomials are always irreducible (see Exercise 9).
Lemma 2.7. If the primitive nth root of unity $\zeta_{n}$ is constructible, then $n$ has the form given in the Gauss-Wantzel Theorem. Moreover, for every $j \geq 1$, each primitive $2^{j}$ th root of unity $\zeta_{2^{j}}$ is constructible.

Proof. Suppose $\zeta_{n}$ is constructible, and let $K_{0} \subset \cdots \subset K_{m}$ be a tower of fields as in Lemma 2.6 ending with $\zeta_{n} \in K_{m}$. Then the irreducibility of the cyclotomic polynomial $\Phi_{n}(T)$ implies

$$
\left[\mathbf{Q}\left(\zeta_{n}\right): \mathbf{Q}\right]=\varphi(n) \mid\left[K_{m}: \mathbf{Q}\right] .
$$

But

$$
\left[K_{m}: K_{m-1}\right]\left[K_{m-1}: K_{m-2}\right] \cdots\left[K_{1}: K_{0}\right]=2^{r}
$$

for some $r \geq 0$. Hence $\varphi(n)$ is a power of 2 , and it is easy to show (Exercise 2) that this forces $n$ to be of the form described in Theorem 2.1.

The final claim of the lemma follows easily by induction: $1=\zeta_{2^{0}}$ is constructible. If all the $2^{j-1}$ th primitive roots of unity are constructible, then so is an arbitrary primitive $2^{j}$ th root of unity $\zeta_{2^{j}}$, since $\left(\zeta_{2^{j}}\right)^{2}$ is primitive of order $2^{j-1}$.

We can reduce the remaining portion of the Gauss-Wantzel result to the following theorem:
Theorem 2.8 (Gauss). Let $p$ be a Fermat prime, and let $\zeta_{p}$ be a primitive $p$ th root of unity. Then $\zeta_{p}$ is constructible.

Suppose Theorem 2.8 is proven. Let $n:=2^{e} p_{1} \cdots p_{k}$ be as in the theorem statement. Since the constructible numbers form a field, it follows that $\zeta_{2^{e}} \zeta_{p_{1}} \cdots \zeta_{p_{r}}$ is constructible (for any choices of the primitive roots of unity in question). But $\zeta_{2} \zeta_{p_{1}} \cdots \zeta_{p_{r}}$ is a primitive $n$th root of unity, and as remarked above, the constructibility of a primitive $n$th root of unity implies the constructibility of the regular $n$-gon.

Below we will give a proof of Theorem 2.8 in the spirit of Gauss. For this it is first necessary to investigate the arithmetic of $\mathbf{Z}\left[\zeta_{p}\right]$.

## 3. Much ado about $\mathrm{Z}\left[\zeta_{p}\right]$

Let $p$ be a prime number, and let $\zeta=\zeta_{p}$ be a complex primitive $p$ th root of unity. In this section we study the arithmetic of $\mathbf{Z}[\zeta]$. Since $\mathbf{Z}[\zeta]$ is the ring of algebraic integers of the cyclotomic field $\mathbf{Q}(\zeta)$, much of this material will be old hat to those versed in algebraic number theory; however, our needs are simple, and we can develop everything that we need from scratch.

Lemma 2.9 (Determination of an integral basis). Every element of $\mathbf{Z}[\zeta]$ (respectively $\mathbf{Q}(\zeta)$ ) can be expressed uniquely in the form $a_{1} \zeta+a_{2} \zeta^{2}+\cdots+$ $a_{p-1} \zeta^{p-1}$, with integral (respectively rational) coefficients $a_{i}$.

Proof. We prove the claim for $\mathbf{Z}[\zeta]$; the proof for $\mathbf{Q}(\zeta)$ is similar. (Note that $\mathbf{Q}(\zeta)=\mathbf{Q}[\zeta]$, since $\zeta$ is algebraic.)
Existence: Since $\zeta$ is a primitive $p$ th root of unity, it is a root of the cyclotomic polynomial

$$
\Phi_{p}(T):=\frac{T^{p}-1}{T-1}=T^{p-1}+T^{p-2}+\cdots+T+1
$$

Substituting $\zeta$ for $T$ yields

$$
\begin{equation*}
\zeta^{p-1}=-1-\zeta-\zeta^{2}-\cdots-\zeta^{p-2} \tag{2.1}
\end{equation*}
$$

This relation together with induction implies that every power of $\zeta$ can be represented as a $\mathbf{Z}$-linear combination of $1, \zeta, \ldots, \zeta^{p-2}$. It then follows that each element of $\mathbf{Z}[\zeta]$ also has a representation of this form. By (2.1), we can write 1 as an integral linear combination of $\zeta, \zeta^{2}, \ldots, \zeta^{p-1}$, and the existence half of Lemma 2.9 follows.
Uniqueness (cf. [Gau86, Art. 341, end of Art. 346]): This is a consequence of the irreducibility of $\Phi_{p}(T)$, which in turn follows from the EisensteinSchönemann criterion:

$$
\Phi_{p}(T+1)=\frac{1}{T}\left((T+1)^{p}-1\right)=\sum_{k=0}^{p-1}\binom{p}{k+1} T^{k}
$$

is a monic polynomial all of whose nonleading coefficients are divisible by $p$, and whose constant coefficient is equal to $p$. Hence $1, \zeta, \zeta^{2}, \ldots, \zeta^{p-2}$ are Q-linearly independent, and so are $\zeta \cdot 1, \zeta \cdot \zeta, \ldots, \zeta \cdot \zeta^{p-2}$.

Remark. See [Gau86, Art. 341] for Gauss's original proof of the irreducibility of $\Phi_{p}(T)$, which was considerably more complicated. In Exercise 9 we show that $\Phi_{n}(T)$ is irreducible for every $n$.

Lemma 2.10. Suppose $\alpha \in \mathbf{Z}[\zeta] \cap \mathbf{Q}$. Then $\alpha \in \mathbf{Z}$. That is, the only rational elements of $\mathbf{Z}[\zeta]$ are the rational integers.

Proof. By Lemma 2.9, we can write $\alpha=a_{1} \zeta+\cdots+a_{p-1} \zeta^{p-1}$ for integers $a_{i}$. Since $\alpha \in \mathbf{Q}$, the expression $\alpha=-\sum_{i=1}^{p-1} \alpha \zeta^{i}$ is a representation of $\alpha$ as a $\mathbf{Q}$-linear combination of $\zeta, \zeta^{2}, \ldots, \zeta^{p-1}$. By the uniqueness half of Lemma 2.9, it follows that $a_{i}=-\alpha$ for each $i$. In particular, $\alpha=-a_{1} \in \mathbf{Z}$.

We turn next to a study of the Galois theory of $\mathbf{Q}(\zeta) / \mathbf{Q}$ :

Lemma 2.11 (Description of the automorphisms of $\mathbf{Q}(\zeta) / \mathbf{Q})$. For each element $a \bmod p \in(\mathbf{Z} / p \mathbf{Z})^{\times}$, there is an automorphism $\sigma_{a}$ of $\mathbf{Q}(\zeta) / \mathbf{Q}$ sending $\zeta \mapsto \zeta^{a}$. Moreover, every such automorphism is of this form. Consequently, $\operatorname{Gal}(\mathbf{Q}(\zeta) / \mathbf{Q})$ can be identified with $(\mathbf{Z} / p \mathbf{Z})^{\times}$.

Proof. The automorphisms of $\mathbf{Q}(\zeta)$ are determined by where they send $\zeta$. The possible images are the roots of $\Phi_{p}$, which are precisely $\zeta^{a}$ for $(a, p)=$ 1. So for each $(a, p)=1$, there is an automorphism $\sigma_{a}$ with $\zeta \mapsto \zeta^{a}$, and these exhaust the automorphisms. Moreover, $\sigma_{a}=\sigma_{a^{\prime}}$ precisely when $a \equiv a^{\prime}(\bmod p)$. Finally, notice that

$$
\sigma_{a} \circ \sigma_{a^{\prime}}(\zeta)=\sigma_{a}\left(\zeta^{a^{\prime}}\right)=\zeta^{a a^{\prime}}=\sigma_{a a^{\prime}}(\zeta)
$$

Putting everything together, we see that the map $a \bmod p \mapsto \sigma_{a}$ is an isomorphism between $(\mathbf{Z} / p \mathbf{Z})^{\times}$and the Galois group $\operatorname{Gal}(\mathbf{Q}(\zeta) / \mathbf{Q})$.

Lemma 2.12 (Description of the fixed fields; cf. [Gau86, Art. 347]). Let $H$ be a subgroup of $(\mathbf{Z} / p \mathbf{Z})^{\times}$; then $H$ is the set of eth powers for a uniquely defined natural number e dividing $p-1$. Write $p-1=e f$.

Let $g$ be a fixed generator of $(\mathbf{Z} / p \mathbf{Z})^{\times}$. Then the set of elements of $\mathbf{Q}(\zeta)$ (respectively $\mathbf{Z}[\zeta]$ ) fixed by $\sigma_{a}$ for every $a \in H$ is precisely the set of $\mathbf{Q}$-linear (resp. Z-linear) combinations of $\eta_{1}, \ldots, \eta_{e}$, where

$$
\begin{equation*}
\eta_{i}:=\zeta^{g^{i}}+\zeta^{g^{e+i}}+\zeta^{g^{2 e+i}}+\cdots+\zeta^{g^{e(f-1)+i}}=\sum_{m=0}^{f-1} \zeta^{g^{e m+i}} \tag{2.2}
\end{equation*}
$$

Following Gauss, we refer to the numbers $\eta_{1}, \ldots, \eta_{e}$ as the $f$-nomial periods (associated to this prime $p$ and this choice of a generator $g$ ). Note that the complex numbers $\eta_{1}, \ldots, \eta_{e}$ are distinct because of Lemma 2.9. It is convenient to take (2.2) as defining $\eta_{i}$ for every integer $i$; then the $\eta_{i}$ are periodic in $i$ with minimal period $e$.

Proof. The assertion that $H$ is the set of $e$ th powers for a unique positive divisor $e$ of $p-1$ follows from the cyclic nature of $(\mathbf{Z} / p \mathbf{Z})^{\times}$. Since $g$ is a generator of $(\mathbf{Z} / p \mathbf{Z})^{\times}$, we have $H=\left\langle g^{e}\right\rangle$. Thus an element of $\mathbf{Q}(\zeta)$ is fixed by everything in $H$ once it is fixed by the single automorphism $\sigma_{g^{e}}$.

Suppose $\alpha$ is fixed by $\sigma_{g^{e}}$. Write $\alpha=\sum_{i=1}^{p-1} c_{i} \zeta^{g^{i}}$, and extend the indices on the $c_{i}$ cyclically with period $p-1$ (i.e., set $c_{i}:=c_{i \bmod p-1}$ for all $i$ ). Lemma 2.9 implies that $\alpha$ is fixed by $\sigma_{g^{e}}$ if and only if $c_{i}=c_{i+e}$ for all $i$. But then

$$
\begin{aligned}
\alpha=c_{1}\left(\zeta^{g}\right. & \left.+\zeta^{g^{e+1}}+\cdots+\zeta^{g^{(f-1) e+1}}\right)+c_{2}\left(\zeta^{g^{2}}+\zeta^{g^{e+2}}+\cdots+\zeta^{g^{(f-1) e+2}}\right) \\
& +\cdots+c_{e}\left(\zeta^{g^{e}}+\zeta^{g^{2 e}}+\cdots+\zeta^{g^{e f}}\right)=c_{1} \eta_{1}+c_{2} \eta_{2}+\cdots+c_{e} \eta_{e}
\end{aligned}
$$

is a linear combination of the $\eta_{i}$, as claimed.

The converse is clear, since each of the $\eta_{i}$ is fixed by $\sigma_{g^{e}}$.
Corollary 2.13. Let $\alpha$ be an element of $\mathbf{Z}[\zeta]$ and suppose that $\sigma_{a}(\alpha)=\alpha$ for every $a \in(\mathbf{Z} / p \mathbf{Z})^{\times}$. Then $\alpha$ is a rational integer.

Proof. We apply the lemma with $H=(\mathbf{Z} / p \mathbf{Z})^{\times}$(and hence $\left.e=1, f=p-1\right)$ to obtain that $\alpha$ is a $\mathbf{Z}$-linear combination of the $(p-1)$-nomial period

$$
\eta_{1}=\sum_{m=0}^{p-2} \zeta^{g^{m+1}}=\zeta+\zeta^{2}+\cdots+\zeta^{p-1}=-1
$$

## 4. Completion of the proof of the Gauss-Wantzel theorem

Suppose that $p$ is a Fermat prime, so that $p-1=2^{n}$ for some positive integer $n$. Let $g$ be a fixed generator of $(\mathbf{Z} / p \mathbf{Z})^{\times}$, and write down the $2^{n}$-nomial period

$$
\begin{equation*}
\zeta^{g^{0}}+\zeta^{g^{1}}+\cdots+\zeta^{g^{p-2}} \tag{2.3}
\end{equation*}
$$

We split this into two $2^{n-1}$-nomial periods by taking every other term,

$$
\begin{equation*}
\zeta^{g^{0}}+\zeta^{g^{2}}+\zeta^{g^{4}}+\cdots+\zeta^{g^{p-1}}, \quad \zeta^{g^{1}}+\zeta^{g^{3}}+\zeta^{g^{5}}+\cdots+\zeta^{g^{p-2}} \tag{2.4}
\end{equation*}
$$

Each of these then splits into two $2^{n-2}$-nomial periods in the same manner. Continuing in this way we eventually reach a level with $2^{n} 1$-nomial periods (which are simply the individual $2^{n}$ primitive $p$ th roots of unity).

To codify this process, we let $\left(2^{n}, g^{0}\right)$ denote the $2^{n}$-nomial period (2.3), we let $\left(2^{n-1}, g^{0}\right)$ and $\left(2^{n-1}, g^{1}\right)$ denote the first and second $2^{n-1}$-nomial periods indicated in (2.4), and in general we let $(f, j)$ denote the $f$-nomial period containing $\zeta^{j}$.

Splitting up the period (2.3) like this yields a binary tree whose first few rows are shown in the following diagram. Here each period is the sum of the two periods from the nodes immediately below:


In general, $\left(2^{n-r}, g^{k}\right)$ branches off (if $\left.r<n\right)$ to yield the two periods $\left(2^{n-r-1}, g^{k}\right)$ and $\left(2^{n-r-1}, g^{k+2^{r}}\right)$. Moreover, the $2^{r}$ periods of the $r$ th row (numbered starting with $r=0$ ) are a complete list of the $2^{n-r}$-nomial periods. To see this, let $f=2^{n-r}$. Then there are $e=(p-1) / f=2^{r}$ distinct


Figure 1. Gauss [Gau86, Art. 354]: Binary tree illustrating (for $p=$ $17, g=3$ ) the decomposition of the 16 -nomial period $\zeta^{1}+\zeta^{3}+\zeta^{9}+$ $\zeta^{10}+\cdots+\zeta^{2}+\zeta^{6}$ into successive half-periods. The correctness of this diagram can be verified with the aid of the following table of powers of $3(\bmod 17)$ :

| $n$ | 0 | 1 | 2 | 3 | 4 | 5 | 6 | 7 | 8 | 9 | 10 | 11 | 12 | 13 | 14 | 15 |
| ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: | ---: |
| $3^{n} \bmod 17$ | 1 | 3 | 9 | 10 | 13 | 5 | 15 | 11 | 16 | 14 | 8 | 7 | 4 | 12 | 2 | 6 |

periods. But the $r$ th row contains $2^{r}$ distinct $2^{n-r}$-nomial periods by construction (each constructed period is distinct from the others by Lemma 2.9). The claim follows.

We can now prove Gauss's result that $\zeta_{p}$ is constructible. By the remarks in the introduction, this will complete the proof of Theorem 2.1.

Proof of Theorem 2.8. Certainly the (unique) $2^{n}$-nomial period is constructible, being just $\zeta+\cdots+\zeta^{p-2}+\zeta^{p-1}=-1$.

Suppose now that every period in the $r$ th row (i.e., every $2^{n-r}$-nomial period) is constructible, for a certain $0 \leq r<n$. Choose a node in the $r$ th row, say $\left(2^{n-r}, g^{k}\right)$, and consider the polynomial

$$
\psi_{r}(T):=\left(T-\left(2^{n-(r+1)}, g^{k}\right)\right)\left(T-\left(2^{n-(r+1)}, g^{k+2^{r}}\right)\right)
$$

whose roots are the periods beneath this node. Since $\sigma_{g^{2^{r}}}\left(\left(2^{n-(r+1)}, g^{k}\right)\right)=$ $\left(2^{n-(r+1)}, g^{k+2^{r}}\right)$ and

$$
\sigma_{g^{2^{r}}}\left(\left(2^{n-(r+1)}, g^{k+2^{r}}\right)\right)=\left(2^{n-(r+1)}, g^{k+2^{r+1}}\right)=\left(2^{n-(r+1)}, g^{k}\right),
$$

the automorphism $\sigma_{g^{2 r}}$ permutes the factors of $\psi_{r}(T)$, and so leaves the coefficients of $\psi_{r}$ fixed. It follows from Lemma 2.12 (with $e=2^{r}, f=$ $2^{n-r}$ ) that the coefficients of $\psi_{r}$ are $\mathbf{Z}$-linear combinations of the $2^{n-r}$-nomial periods. In particular, they are constructible by the induction hypothesis.

Since the constructible numbers form a field closed under the taking of square roots, the quadratic formula shows that both roots $\left(2^{n-(r+1)}, g^{k}\right)$ and $\left(2^{n-(r+1)}, g^{k+2^{r}}\right)$ of $\psi_{r}$ are constructible.

Proceeding like this for each node in the $r$ th row, we obtain the constructibility of all the periods in the $(r+1)$ th row. Theorem 2.8 now follows by induction, since the individual primitive $p$ th roots of unity are the $\left(2^{0}-\right.$ nomial) periods of the $n$th row.

A detailed treatment of the case $p=17$ is the subject of Exercise 6 .

## 5. Period polynomials and Kummer's criterion

If $p \equiv 1(\bmod e)$ is prime, then the period polynomial $\phi(T)$ of degree $e$ is defined by

$$
\phi(T):=\left(T-\eta_{0}\right)\left(T-\eta_{1}\right) \cdots\left(T-\eta_{e-1}\right) \in \mathbf{C}[T],
$$

and the reduced period polynomial $\hat{\phi}(x)$ of degree $e$ is defined by

$$
\hat{\phi}(T):=\left(T-\left(e \eta_{0}+1\right)\right)\left(T-\left(e \eta_{1}+1\right)\right) \cdots\left(T-\left(e \eta_{e-1}+1\right)\right),
$$

where the $\eta_{i}$ are the $f$-nomial periods (and, as usual, $p=e f+1$ ). Note that since the choice of a generator $g$ of $(\mathbf{Z} / p \mathbf{Z})^{\times}$only impacts the order of the $\eta_{i}$, both $\phi$ and $\hat{\phi}$ are independent of the choice of $g$.

At this point $\phi$ is arguably as natural to introduce as the Gaussian periods themselves. But what is $\hat{\phi}$ ? We can describe $\hat{\phi}$ by describing its roots: They are

$$
\begin{equation*}
e \eta_{0}+1=1+e \sum_{m=0}^{f-1} \zeta^{g^{e m}}=1+e \sum_{u \bmod p \in\left(\mathbf{F}_{p}^{\times}\right)^{e}} \zeta^{e}=\sum_{v \bmod p \in \mathbf{F}_{p}} \zeta^{v^{e}} . \tag{2.5}
\end{equation*}
$$

and its images under the various automorphisms $\sigma_{a}$. For us, the importance of $\hat{\phi}$ rests in the observation that

$$
\begin{aligned}
\sum_{i}\left(e \eta_{i}+1\right) & =e \sum_{i} \eta_{i}+e \\
& =e\left(1+\zeta+\cdots+\zeta^{p-1}\right)+e=-e+e=0
\end{aligned}
$$

so that the next-to-leading coefficient of $\hat{\phi}$ automatically vanishes. This makes $\hat{\phi}$ a simpler object to work with.

We now prove that $\phi$ and $\hat{\phi}$, which a priori have complex coefficients, in fact have integer coefficients and are irreducible over the rationals:

Theorem 2.14. The period polynomial $\phi(T)$ has integer coefficients and is irreducible over the rationals. The same holds for $\hat{\phi}$.

Of course this agrees with what we already know about the $p$ th cyclotomic polynomial (which corresponds to taking $e=p-1, f=1$ ). Below we will compute the period polynomials and reduced period polynomials of degree 2 and 3 .

Proof of Theorem 2.14. It suffices to prove only the statements for $\phi$ owing to the relation

$$
\begin{equation*}
\hat{\phi}(T)=\prod_{i=0}^{e-1}\left(T-\left(e \eta_{i}+1\right)\right)=e^{e} \prod_{i=0}^{e-1}\left(\frac{T-1}{e}-\eta_{i}\right)=e^{e} \phi((T-1) / e) . \tag{2.6}
\end{equation*}
$$

The coefficients of $\phi(T)$ belong to $\mathbf{Z}[\zeta]$, so (by Corollary 2.13) to show that they are rational integers, it is enough to check that they are fixed by every $\sigma_{a}$. Assume that the $\eta_{i}$ are defined with respect to the generator $g$ of $\mathbf{F}_{p}^{\times}$. If the index of $a$ with respect to $g$ is congruent to $i(\bmod e)$, then $\sigma_{a}\left(\eta_{j}\right)=\eta_{i+j}$. Since $i+j$ runs through a complete residue system modulo $e$ as $j$ does, it follows that $\sigma_{a}$ merely permutes the roots of $\phi(T)$, and so fixes its coefficients.

Irreducibility is surprisingly easy: Given a polynomial over the rationals which vanishes at $\eta_{0}$, we repeatedly apply the automorphism $\sigma_{g}$ to see that this polynomial also vanishes at $\eta_{1}, \eta_{2}, \ldots$. Since the $\eta_{i}$ are distinct, the given polynomial must be divisible by $\phi$. This implies that $\phi$ generates the
ideal of polynomials in $\mathbf{Q}[T]$ which vanish at $\eta_{0}$. This is a prime ideal, hence $\phi$ itself is prime.

The next theorem provides the link between period polynomials and the study of higher reciprocity. Keeping with tradition, we have attributed it to Kummer (see [Kum46]), but it appears to have been known earlier to Gauss (cf. [Gau65, Art. 367]):

Theorem 2.15 (Kummer's criterion). Let $p=e f+1$ be prime, and let $\phi$ be the period polynomial of degree $e$. Let $q$ be a prime distinct from $p$.
(i) If $q$ is an eth power modulo $p$, then the polynomial $\phi(T)$ has a root mod $q$.
(ii) Conversely, if $q$ is a prime not dividing the discriminant of $\phi$ for which $\phi$ has a root $\bmod q$, then $q$ is an eth power residue $\bmod p$.
(iii) Suppose moreover that $e$ is prime. Then every $q$ dividing the discriminant of $\phi$ is an eth power residue of $p$.

When $e$ is prime, statements (i)-(iii) have the following elegant corollary:
Corollary 2.16. With notation as in Theorem 2.15, q is an eth power residue modulo $p$ if and only if $\phi$ has a root modulo $q$.

The proof of Theorem 2.15 requires the following simple lemma:
Lemma 2.17. Keep the notation of Theorem 2.15. Suppose that $\eta_{i} \equiv$ $\eta_{j}(\bmod q)$, where the congruence is in the ring $\mathbf{Z}[\zeta]$. Then $i \equiv j(\bmod e)$.

Proof. If $\eta_{i} \equiv \eta_{j}(\bmod q)$, then $q$ divides $\eta_{i}-\eta_{j}$. Lemma 2.9 then implies that $q$ divides every coefficient of $\eta_{i}-\eta_{j}$ when both are expressed as Z-linear combinations of $\zeta, \zeta^{2}, \ldots, \zeta^{p-1}$. But referring to the definition (2.2) of the $\eta_{i}$ shows that this is only possible when $\eta_{i}=\eta_{j}$, i.e., when $i \equiv j(\bmod e)$.

Proof of Theorem 2.15. We work modulo $q$ in the ring $\mathbf{Z}[\zeta]$. Fix a generator $g$ of $(\mathbf{Z} / p \mathbf{Z})^{\times}$, and use this generator to determine the numbering of the periods $\eta_{i}$. Suppose that $q \equiv g^{r}(\bmod p)$. From the binomial theorem,

$$
\eta_{k}^{q}=\left(\sum_{m=0}^{f-1} \zeta^{g^{e m+k}}\right)^{q} \equiv \sum_{m=0}^{f-1} \zeta^{g^{e m+k+r}} \equiv \eta_{k+r}(\bmod q)
$$

Now let $n$ be an arbitrary integer. Since $y^{q}-y=\prod_{i=0}^{q-1}(y-i)$ is an identity in every ring of characteristic $q$, we have

$$
\begin{aligned}
\left(n-\eta_{k}\right)\left(n-\eta_{k}-1\right) \cdots\left(n-\eta_{k}-(q-1)\right) & \equiv\left(n-\eta_{k}\right)^{q}-\left(n-\eta_{k}\right) \\
& \equiv \eta_{k}-\eta_{k}^{q} \equiv \eta_{k}-\eta_{k+r}(\bmod q) .
\end{aligned}
$$

Multiplying over $k=0,1, \ldots, e-1$, we obtain

$$
\begin{equation*}
\phi(n) \phi(n-1) \cdots \phi(n-(q-1)) \equiv \prod_{k=0}^{e-1}\left(\eta_{k}-\eta_{k+r}\right)(\bmod q) \tag{2.7}
\end{equation*}
$$

If $q$ is an $e$ th power modulo $p$, then $e$ divides $r$, and so $\eta_{k+r}=\eta_{k}$ for each $k$. Hence $q$ divides $\phi(n) \cdots \phi(n-(q-1))$ in $\mathbf{Z}[\zeta]$. By Lemma 2.10, the same divisibility relation holds over $\mathbf{Z}$. Since $q$ is prime in $\mathbf{Z}$, it follows that $q$ divides (over the integers) some value of $\phi$, which is the assertion of (i).

The congruence (2.7) also yields a quick proof of (ii): If $q \mid \phi(n)$ and $q$ is not an $e$ th power residue $\bmod p$, then $e \nmid r$. Hence, defining

$$
P_{j}:=\prod_{k=0}^{e-1}\left(\eta_{k}-\eta_{k+j}\right), \quad \text { we have } \quad q\left|P_{r}\right| \prod_{j=1}^{e-1} P_{j} \mid \operatorname{Disc}(\phi)
$$

in $\mathbf{Z}[\zeta]$. The same divisibility holds also in $\mathbf{Z}$, and this proves (ii).
We now prove (iii). We suppose that $q$ divides the discriminant of $\phi$ and show that in this case $e \mid r$, so that $q \equiv g^{r}$ must be an eth power residue.

Suppose instead that $e \nmid r$. Then $r$ is coprime to $e$, since $e$ is a rational prime by hypothesis. Now the $P_{j}$ are rational integers, since they are fixed by every automorphism $\sigma_{a}$. Since

$$
q \mid \operatorname{Disc}(\phi)= \pm \prod_{1 \leq j \leq e-1} P_{j}
$$

we can choose an index $j, 1 \leq j \leq e-1$, for which $q \mid P_{j}$. Then

$$
\begin{aligned}
\left(\eta_{0}-\eta_{j}\right)^{\frac{e^{e}-1}{q-1}}=\prod_{i=0}^{e-1}\left(\eta_{0}-\eta_{j}\right)^{q^{i}} & \equiv \prod_{i=0}^{e-1}\left(\eta_{i r}-\eta_{i r+j}\right) \\
& \equiv \prod_{i=0}^{e-1}\left(\eta_{i}-\eta_{i+j}\right) \equiv P_{j}(\bmod q)
\end{aligned}
$$

using that $r$ is coprime to $e$, so that $i r$ runs through a complete residue system modulo $e$ as $i$ does. Since $q \mid P_{j}$, it follows that

$$
q\left|\left(\eta_{0}-\eta_{j}\right)^{\frac{q^{e}-1}{q-1}}\right|\left(\eta_{0}-\eta_{j}\right)^{q^{e}},
$$

and so

$$
0 \equiv\left(\eta_{0}-\eta_{j}\right)^{q^{e}} \equiv \eta_{0+r e}-\eta_{j+r e} \equiv \eta_{0}-\eta_{j}(\bmod q)
$$

so that $\eta_{0} \equiv \eta_{j}(\bmod e)$. But this contradicts Lemma 2.17.

## 6. A cyclotomic proof of quadratic reciprocity

Let $p$ be an odd prime. Then $p-1$ is even, and so it makes sense to consider the period polynomial of degree 2 . We will prove quadratic reciprocity by applying Kummer's criterion (Theorem 2.15) with $e=2$. For this we need an explicit determination of the quadratic period polynomial:

Theorem 2.18. Let $p$ be an odd prime, and put $p^{*}=(-1)^{(p-1) / 2} p$, so that $p^{*}=p$ if $p \equiv 1(\bmod 4)$ and $p^{*}=-p$ otherwise. The period polynomial of degree $e=2$ is

$$
T^{2}+T+\frac{1-p^{*}}{4}
$$

The reduced period polynomial of degree 2 is $T^{2}-p^{*}$.

The proof of this theorem will be facilitated by means of the following lemma, which allows us to simplify any product of two $f$-nomial periods (where, as usual, we write $p=e f+1$ ). Before we can state the lemma, we need to introduce the cyclotomic numbers. Fix a generator $g$ of $\mathbf{F}_{p}^{\times}$. If $\alpha \in \mathbf{F}_{p}^{\times}$, the index of $\alpha$ (with respect to $g$ ), denoted $\operatorname{ind}_{g} \alpha$, is the integer $k \in[0, p-2]$ for which $g^{k}=\alpha$. The cyclotomic numbers are defined for every pair of integers $i$ and $j$ by

$$
\begin{equation*}
(i, j):=\sum_{\substack{\alpha \in \mathbf{F}_{p} \backslash\{0,-1\} \\ \operatorname{ind} g \alpha=i(\bmod e) \\ \operatorname{ind}_{g}(\alpha+1) \equiv j(\bmod e)}} 1 . \tag{2.8}
\end{equation*}
$$

While we have made this definition for all pairs of $i$ and $j$, of course $i$ and $j$ really only matter modulo $e$. (In our contexts there will be no danger of confusing this " $(i, j)$ " with that used to identify the periods of Fermat primes previously.)

Lemma 2.19. Let $p \equiv 1(\bmod e)$ be prime, and write $p=e f+1$. Let $\eta_{1}, \ldots, \eta_{e}$ denote the $f$-nomial periods. We assume that both the $f$-nomial periods and the cyclotomic numbers are indexed with respect to the same fixed generator $g \bmod p$ of $(\mathbf{Z} / p \mathbf{Z})^{\times}$. Then for every pair of integers $i$ and j, we have

$$
\eta_{i} \eta_{i+j}=\sum_{m=0}^{e-1}(j, m) \eta_{i+m}+ \begin{cases}f & \text { if } j \equiv e f / 2(\bmod e)  \tag{2.9}\\ 0 & \text { otherwise }\end{cases}
$$

Proof. We have

$$
\begin{aligned}
\eta_{i} \eta_{i+j}=\sum_{m=0}^{f-1} \zeta^{g^{e m+i}} & \sum_{n=0}^{f-1} \zeta^{g^{e n+i+j}}=\sum_{m=0}^{f-1} \sum_{n=0}^{f-1} \zeta^{g^{e m+i}}\left(1+g^{e(n-m)+j}\right) \\
& \left.=\sum_{m=0}^{f-1} \sum_{n=0}^{f-1} \zeta^{g^{e m+i}\left(1+g^{e n+j}\right.}\right)=\sum_{n=0}^{f-1} \sum_{m=0}^{f-1} \zeta^{g^{e m+i}}\left(1+g^{e n+j}\right)
\end{aligned}
$$

where in the transition from the first line to the second we use that $n-m$ runs over a complete residue system modulo $f$ as $n$ does (for fixed $m$ ).

Suppose $n$ is such that $\operatorname{ind}_{g}\left(1+g^{e n+j}\right) \equiv r(\bmod e)$. Then the inner sum over $m$ (for this $n$ ) is $\eta_{i+r}$. The number of values of $n$ with $0 \leq n \leq f-1$ for which $\operatorname{ind}_{g}\left(1+g^{e n+j}\right) \equiv r(\bmod e)$ is the cyclotomic number $(j, r)$. Adding the contributions from $r=0,1, \ldots, e-1$ gives the main term in (2.9).

The secondary term comes from the (unique if it exists) value of $n$ with $0 \leq n \leq f-1$ for which $1+g^{e n+j} \equiv 0(\bmod p)$; this term appears if and only if $(p-1) / 2=e f / 2 \equiv j(\bmod e)$.

Proof of Theorem 2.18. We have

$$
\phi(T)=\left(T-\eta_{0}\right)\left(T-\eta_{1}\right)=T^{2}-\left(\eta_{0}+\eta_{1}\right) T+\eta_{0} \eta_{1}
$$

We have

$$
\begin{aligned}
\eta_{0}+\eta_{1} & =\left(\zeta^{g^{0}}+\zeta^{g^{2}}+\cdots+\zeta^{g^{p-1}}\right)+\left(\zeta^{g^{1}}+\zeta^{g^{3}}+\cdots+\zeta^{g^{p-2}}\right) \\
& =\sum_{a \bmod p \in(\mathbf{Z} / p \mathbf{Z})^{\times}} \zeta^{a}=-\zeta^{0}=-1
\end{aligned}
$$

and it remains only to compute $\eta_{0} \eta_{1}$. By Lemma 2.19 with $e=2$ and $f=(p-1) / 2$, we have

$$
\eta_{0} \eta_{1}=(1,0) \eta_{0}+(1,1) \eta_{1}+ \begin{cases}f & \text { if } f \text { is odd } \\ 0 & \text { if } f \text { is even }\end{cases}
$$

The automorphism $\sigma_{g}$ interchanges $\eta_{0}$ and $\eta_{1}$ and hence leaves $\eta_{0} \eta_{1}$ fixed. From the expression just obtained for $\eta_{0} \eta_{1}$ and the $\mathbf{Q}$-linear independence of $\eta_{0}$ and $\eta_{1}$ (coming from Lemma 2.9), we must have $(1,0)=(1,1)$. Hence

$$
\begin{aligned}
2(1,1)=(1,1)+(1,0) & =\sum_{\substack{\alpha \in \mathbf{F}_{p} \backslash\{0,-1\} \\
\operatorname{ind}_{g} \alpha=1(\bmod 2)}} 1 \\
& =\sum_{\substack{1 \leq a<p-1 \\
\left(\frac{a}{p}\right)=-1}} 1=\frac{p-1}{2}-\frac{1-\left(\frac{-1}{p}\right)}{2},
\end{aligned}
$$

If $p \equiv 1(\bmod 4)$, then $f$ is even and $\left(\frac{-1}{p}\right)=1$. Hence

$$
\begin{aligned}
\eta_{0} \eta_{1}=(1,0) \eta_{0}+(1,1) \eta_{1} & =\left(\eta_{0}+\eta_{1}\right)(1,1)=-(1,1) \\
& =-\frac{1}{2}\left(\frac{p-1}{2}\right)=\frac{1-p}{4}=\frac{1-p^{*}}{4} .
\end{aligned}
$$

If $p \equiv 3(\bmod 4)$, then $f$ is odd and $\left(\frac{-1}{p}\right)=-1$, so that

$$
\begin{aligned}
& \eta_{0} \eta_{1}=(1,0) \eta_{0}+(1,1) \eta_{1}+\frac{p-1}{2}=-(1,1)+\frac{p-1}{2} \\
&=-\frac{1}{2}\left(\frac{p-3}{2}\right)+\frac{p-1}{2}=\frac{1+p}{4}=\frac{1-p^{*}}{4}
\end{aligned}
$$

This proves the claim about the form of the period polynomial. It follows from (2.6) that the reduced period polynomial is $4 \phi(T / 2-1 / 2)=T^{2}-p^{*}$, which finishes the proof.

We are now almost in a position to prove quadratic reciprocity. The only additional ingredient required is the following basic result:

Lemma 2.20 (First supplementary law). For each odd prime $p$, we have $\left(\frac{-1}{p}\right)=(-1)^{(p-1) / 2}$.

The proof is trivial: A square root of -1 exists modulo $p$ exactly when there is a primitive fourth root of unity in $\mathbf{F}_{p}^{\times}$. Since $\mathbf{F}_{p}^{\times}$is cyclic, the latter occurs exactly when $p \equiv 1(\bmod 4)$. It is easy to check that this agrees with the answer provided by Lemma 2.20.

Proof of quadratic reciprocity (Theorem 2.2). Let $p$ and $q$ be distinct odd primes. Then $q$ does not divide the discriminant $p^{*}$ of the period polynomial $T^{2}+T+\frac{1}{4}\left(1-p^{*}\right)$. By parts (i) and (ii) of Kummer's criterion (Theorem 2.15),

$$
\begin{aligned}
\left(\frac{q}{p}\right)=1 & \Longleftrightarrow T^{2}+T+\frac{1-p^{*}}{4} \text { has a root modulo } q \\
& \Longleftrightarrow \operatorname{Disc}\left(T^{2}+T+\frac{1-p^{*}}{4}\right) \text { is a square } \bmod q \Longleftrightarrow\left(\frac{p^{*}}{q}\right)=1
\end{aligned}
$$

Thus $\left(\frac{q}{p}\right)=\left(\frac{p^{*}}{q}\right)$. From Lemma 2.20 and the multiplicativity of the Legendre symbol, we have

$$
\left(\frac{p^{*}}{q}\right)=\left(\frac{(-1)^{(p-1) / 2} p}{q}\right)=(-1)^{\frac{p-1}{2} \frac{q-1}{2}}\left(\frac{p}{q}\right),
$$

which gives Theorem 2.2.

This proof of quadratic reciprocity most closely resembles the demonstration offered by V.A. Lebesgue [Leb60]. However, the same ideas can already be found in Gauss's third and fourth proofs of quadratic reciprocity [Gau65, Art. 365-366], which were originally intended to be included in the Disquisitiones (see [Fre07]).

Using the same method we can classify the primes for which 2 is a square:
Theorem 2.21 (Second supplementary law). If $p$ is an odd prime, then $\left(\frac{2}{p}\right)=1$ if $p \equiv \pm 1(\bmod 8)$ and $\left(\frac{2}{p}\right)=-1$ if $p \equiv \pm 3(\bmod 8)$.

Proof of the second supplementary law. Let $p$ be an odd prime. Since $2 \nmid p^{*}$, Theorem 2.15 implies that

$$
\begin{aligned}
\left(\frac{2}{p}\right)=1 & \Longleftrightarrow T^{2}+T+\frac{1-p^{*}}{4} \text { has a root } \bmod 2 \\
& \Longleftrightarrow \frac{1-p^{*}}{4} \equiv 0(\bmod 2) \Longleftrightarrow p \equiv \pm 1(\bmod 8) .
\end{aligned}
$$

## 7. Jacobi's cubic reciprocity law

The proof of Jacobi's cubic reciprocity law is entirely analogous to the proof of the quadratic reciprocity law offered in $\S 6$. But each of the corresponding steps is much more difficult; in particular, determining the coefficients of the cubic period polynomial corresponding to a prime $p \equiv 1(\bmod 3)$ requires a considerable amount of ingenuity. Here we follow Gauss's treatment [Gau86, Art. 358] with minor changes in notation. Along the way we will compute the cyclotomic numbers $(i, j)$ of order 3 , which will be used to determine the cubic residue status of 2 and 3 .

Even after we can write down the cubic period polynomial, it is not obvious how to determine whether it has a root modulo a prime $q$; we will tackle this problem by writing down the roots explicitly (in a finite extension of $\mathbf{F}_{q}$ ) using Cardano's formulas and then using properties of the $q$ th power map to detect when a root lies in $\mathbf{F}_{q}$.

### 7.1. Article 358: The cubic period polynomial.

Theorem 2.22 (Determination of the cubic period polynomial). Let $p \equiv$ $1(\bmod 3)$ be prime, say $p=3 f+1$. Write $4 p=L^{2}+27 M^{2}$ with integers $L$ and $M$, where the sign of $L$ is chosen so that $L \equiv 1(\bmod 3)$. Put $L=3 k-2$. Then the cubic period polynomial corresponding to $p$ is

$$
T^{3}+T^{2}-f T-\frac{f+k p}{9}
$$

Theorem 2.23 (Determination of the cyclotomic numbers of order 3). The matrix of cyclotomic numbers

$$
\left(\begin{array}{ccc}
(0,0) & (0,1) & (0,2)  \tag{2.10}\\
(1,0) & (1,1) & (1,2) \\
(2,0) & (2,1) & (2,2)
\end{array}\right) \quad \text { has the shape } \quad\left(\begin{array}{lll}
a & b & c \\
b & c & d \\
c & d & b
\end{array}\right) .
$$

Here $a, b, c$ and $d$ can be described explicitly as follows: we have

$$
a=\frac{f+k}{3}-1 \quad \text { and } \quad d=\frac{f+k}{3} .
$$

We can choose our generator $g$ of $(\mathbf{Z} / p \mathbf{Z})^{\times}$so that either of $b-c=M$ or $b-c=-M$ holds. If $g$ is chosen so that $b-c=M$, then

$$
\begin{equation*}
b=\frac{M}{2}+\frac{2 f-k}{6} \quad \text { and } \quad c=-\frac{M}{2}+\frac{2 f-k}{6} ; \tag{2.11}
\end{equation*}
$$

otherwise these are interchanged.
It appears from Gauss's mathematical diary that he discovered these results on October 1, 1796 [Gra84, Entry 39].

We will prove Theorems 2.22 and 2.23 simultaneously. We first need some easy properties of the cubic cyclotomic numbers:

Lemma 2.24. Let $p \equiv 1(\bmod 3)$ and write $p-1=3 f$. Then the cyclotomic numbers ( $i, j$ ) defined in (2.8) have the following properties:
(i) For every pair of integers $i$ and $j$, we have $(i, j)=(j, i)$.
(ii) We have
(a) $(0,0)+(0,1)+(0,2)=f-1$,
(b) $(1,0)+(1,1)+(1,2)=f$,
(c) $(2,0)+(2,1)+(2,2)=f$.

Proof. Since -1 is a cube in $(\mathbf{Z} / p \mathbf{Z})^{\times}$, the map $\alpha \mapsto-1-\alpha$ is a bijection between the set counted by $(i, j)$ and that counted by $(j, i)$. This proves (i). To prove (ii), note that

$$
(i, 0)+(i, 1)+(i, 2)=\sum_{\substack{\alpha \in \mathbf{F}_{g} \backslash\{0,1\} \\ \operatorname{ind} g_{g} \equiv i(\bmod 3) \\ \operatorname{ind}_{g}(\alpha+1) \equiv 0,1, \text { or } 2(\bmod 3)}} 1
$$

That is, $(i, 0)+(i, 1)+(i, 2)$ counts the number of $\alpha{\text { with } \operatorname{ind}_{g}(\alpha) \equiv i(\bmod 3)}(\alpha)$ and $\alpha+1 \neq 0$. There are $(p-1) / 3=f$ elements $\alpha$ with $\operatorname{ind}_{g}(\alpha) \equiv i(\bmod 3)$. If $i \not \equiv 0(\bmod 3)$, then none of these satisfy $\alpha+1=0$. However, if $i \equiv$ $0(\bmod 3)$, then $\alpha:=-1$ has index congruent to $i \bmod 3$ and $\alpha+1=0$; this explains the anomalous count for $(0,0)+(0,1)+(0,2)$.

Write the period polynomial $\phi(T)$ in the form

$$
\begin{align*}
& T^{3}-A T^{2}+B T-C,  \tag{2.12}\\
& \text { where } \quad A=\eta_{0}+\eta_{1}+\eta_{2}, B=\eta_{0} \eta_{1}+\eta_{1} \eta_{2}+\eta_{0} \eta_{2} \text { and } C=\eta_{0} \eta_{1} \eta_{2}
\end{align*}
$$

are the elementary symmetric functions of $\eta_{0}, \eta_{1}$ and $\eta_{2}$. We have

$$
A=\eta_{0}+\eta_{1}+\eta_{2}=\sum_{a \bmod p \in(\mathbf{Z} / p \mathbf{Z})^{\times}} \zeta^{a}=-1 .
$$

By Lemma 2.19,

$$
\begin{equation*}
\eta_{0} \eta_{1}=(1,0) \eta_{0}+(1,1) \eta_{1}+(1,2) \eta_{2} . \tag{2.13}
\end{equation*}
$$

Applying the automorphism $\sigma_{g}$ we obtain the two further relations

$$
\begin{align*}
& \eta_{1} \eta_{2}=(1,0) \eta_{1}+(1,1) \eta_{2}+(1,2) \eta_{0},  \tag{2.14}\\
& \eta_{2} \eta_{0}=(1,0) \eta_{2}+(1,1) \eta_{0}+(1,2) \eta_{1} . \tag{2.15}
\end{align*}
$$

Adding (2.13), (2.14), and (2.15) we find that

$$
B=\eta_{0} \eta_{1}+\eta_{1} \eta_{2}+\eta_{2} \eta_{3}=((1,0)+(1,1)+(1,2))\left(\eta_{0}+\eta_{1}+\eta_{2}\right)=-f .
$$

Lemma 2.19 also yields

$$
\eta_{0} \eta_{2}=(2,0) \eta_{0}+(2,1) \eta_{1}+(2,2) \eta_{2} .
$$

Comparing this with $(2.15)$, we see that $(2,0)=(1,1)$ and $(2,2)=(1,0)$. This, together with the first statement of Lemma 2.24, proves that the matrix of cyclotomic numbers has the form stated in (2.10). Henceforth we refer to the cyclotomic numbers by their letter designation in that matrix.

By Lemma 2.24,

$$
a+b+c=(0,0)+(0,1)+(0,2)=f-1 \quad \text { and } \quad b+c+d=f,
$$

and so we obtain the additional relation

$$
a=d-1 .
$$

From Lemma 2.19 and equations (2.13), (2.14), (2.15), we have

$$
\begin{array}{ll}
\eta_{0} \eta_{0}=f+(d-1) \eta_{0}+b \eta_{1}+c \eta_{2}, \\
\eta_{0} \eta_{1}= & b \eta_{0}+c \eta_{1}+d \eta_{2}, \\
\eta_{0} \eta_{2}= & c \eta_{0}+d \eta_{1}+b \eta_{2}, \\
\eta_{1} \eta_{2}= & d \eta_{0}+b \eta_{1}+c \eta_{2} .
\end{array}
$$

Hence
(2.16) $=d f+\left(b^{2}+c^{2}+d^{2}-d\right) \eta_{0}+(b d+b c+c d) \eta_{1}+(b d+b c+c d) \eta_{2}$.

Since $C$ is a rational integer, it is fixed by the automorphism $\sigma_{g}$. This automorphism cyclically permutes $\eta_{0}, \eta_{1}$, and $\eta_{2}$, and so the linear independence of the $\eta_{i}$ implies that the coefficients of $\eta_{0}, \eta_{1}$ and $\eta_{2}$ in (2.16) must coincide. That is,

$$
\begin{equation*}
b^{2}+c^{2}+d^{2}-d=b d+b c+c d \tag{2.17}
\end{equation*}
$$

Hence

$$
\begin{aligned}
C & =d f+(b d+b c+c d)\left(\eta_{0}+\eta_{1}+\eta_{2}\right) \\
& =d(b+c+d)-(b d+b c+c d)=d^{2}-b c .
\end{aligned}
$$

Relation (2.17) can also be written in the form

$$
\begin{aligned}
& 12 d+12 b+12 c+4 \\
& \quad=36 d^{2}+36 b^{2}+36 c^{2}-36 b d-36 c d-36 b c-24 d+12 b+12 c+4
\end{aligned}
$$

or, observing that $12(b+c+d)+4=12 f+4=4 p$, very concisely as

$$
4 p=(6 d-3 b-3 c-2)^{2}+27(b-c)^{2} .
$$

(Note that this gives another proof of the existence half of Lemma 2.3.) We began by assuming that $4 p=L^{2}+27 M^{2}$. Since $L$ and $6 d-3 b-3 c-2$ both belong to the residue class $1 \bmod 3$, the uniqueness half of Lemma 2.3 implies that

$$
L=3 k-2=6 d-3 b-3 c-2 \quad \text { and } \quad b-c= \pm M,
$$

so that

$$
k=2 d-b-c=3 d-f .
$$

Hence

$$
\begin{equation*}
d=\frac{f+k}{3} \quad \text { and } \quad b+c=f-d=\frac{2 f-k}{3} . \tag{2.18}
\end{equation*}
$$

Consequently,

$$
\begin{aligned}
C=d^{2}-b c & =d^{2}-\frac{(b+c)^{2}}{4}+\frac{(b-c)^{2}}{4} \\
& =\frac{(f+k)^{2}}{9}-\frac{(2 f-k)^{2}}{36}+\frac{M^{2}}{4} .
\end{aligned}
$$

If we substitute $M^{2}=\frac{1}{27}\left((12 f+4)-(3 k-2)^{2}\right)$, this simplifies to

$$
\frac{k(3 f+1)+f}{9}=\frac{f+k p}{9},
$$

and this finishes the proof of Theorem 2.22.
It is now easy to complete the determination of the cyclotomic numbers. First, replacing $g$ with $g^{-1}$ has the effect of interchanging $b=(0,1)$ and $c=(0,2)$, so that $b-c= \pm M$ can be made to hold for either choice of sign, as was claimed in Theorem 2.23. Next, if $g$ is chosen so that $b-c=M$,
then (2.18) yields (2.11). Similar considerations apply if $g$ is chosen so that $b-c=-M$. This completes the proof of Theorem 2.23.

Corollary 2.25. Let $p \equiv 1(\bmod 3)$. Then $T^{3}-3 p T-p L$ is the reduced cubic period polynomial corresponding to $p$.

Proof. By (2.6) and Theorem 2.22,

$$
\hat{\phi}(T)=3^{3} \phi(T / 3-1 / 3)=T^{3}-3(3 f+1) T+6 f-3 k p+2 .
$$

The corollary follows once we observe that

$$
3 f+1=p \quad \text { and } \quad 6 f-3 k p+2=-3 k p+2 p=p(2-3 k)=-p L .
$$

### 7.2. The cubic character of 2 and 3.

Theorem 2.26 (Cubic character of 2$)$. Let $p \equiv 1(\bmod 3)$, and write $4 p=L^{2}+27 M^{2}$, where $L \equiv 1(\bmod 3)$. Suppose $g$ is a primitive root chosen so that $b-c=M$, where $b=(2,2)$ and $c=(1,1)$ are the cyclotomic numbers of the previous section. Then

2 is a cube $\Longleftrightarrow 2 \mid L$ and $2 \mid M$,

$$
\begin{aligned}
& \operatorname{ind}_{g}(2) \equiv 1(\bmod 3) \Longleftrightarrow 4 \mid L-M \\
& \operatorname{ind}_{g}(2) \equiv 2(\bmod 3) \Longleftrightarrow 4 \mid L+M
\end{aligned}
$$

In particular, 2 is a cube modulo the prime $p \equiv 1(\bmod 3)$ if and only if $p$ can be written in the form $L^{\prime 2}+27 M^{\prime 2}$ for some integers $L^{\prime}$ and $M^{\prime}$.

Proof. Suppose $i \in\{0,1,2\}$. We let $S$ be the set of $\alpha$ counted by the cyclotomic number $(i, i)$. In other words, $S$ is the set of $\alpha \in \mathbf{F}_{p} \backslash\{0,-1\}$ for which $\operatorname{ind}_{g} \alpha \equiv \operatorname{ind}_{g}(\alpha+1) \equiv i(\bmod 3)$. It is easy to check that the map $\psi$ defined on $S$ by $\psi(\alpha)=-1-\alpha$ is an involution of $S$. Since $(i, i)=\# S$,

$$
\begin{align*}
(i, i) \text { is odd } & \Longleftrightarrow \psi \text { has a fixed point } \\
& \Longleftrightarrow \operatorname{ind}_{g}(-1 / 2) \equiv i(\bmod 3)  \tag{2.19}\\
& \Longleftrightarrow \operatorname{ind}_{g}(2) \equiv-i(\bmod 3)
\end{align*}
$$

Since $f=\frac{p-1}{3}$ is even and $L=3 k-2 \equiv-k-2(\bmod 4)$, Theorem 2.23 implies that

$$
\begin{gathered}
(0,0)=a=d-1=\frac{f+k}{3}-1 \equiv k-1 \equiv L-1(\bmod 2) \\
(1,1)=c, \text { and } 2 c=\frac{2 f-k}{3}-M \equiv k-2 f-M \equiv-L-M-2(\bmod 4), \\
(2,2)=b, \text { and } 2 b=M+\frac{2 f-k}{3} \equiv M+k-2 f \equiv M-L-2(\bmod 4) .
\end{gathered}
$$

Theorem 7.5 now follows from the equivalences (2.19): For example, taking $i=0$, we see that

$$
\begin{aligned}
2 \text { is a cube } & \Longleftrightarrow \operatorname{ind}_{g}(2) \equiv 0(\bmod 3) \\
& \Longleftrightarrow(0,0) \text { is odd } \Longleftrightarrow L-1 \text { is odd } \Longleftrightarrow L \text { is even. }
\end{aligned}
$$

The other results are proved similarly:

$$
\begin{aligned}
\operatorname{ind}_{g}(2) \equiv 1(\bmod 3) & \Longleftrightarrow(2,2) \text { is odd } \\
& \Longleftrightarrow 2(2,2) \equiv 2(\bmod 4) \Longleftrightarrow M-L \equiv 0(\bmod 4)
\end{aligned}
$$

and

$$
\begin{aligned}
\operatorname{ind}_{g}(2) \equiv 2(\bmod 3) & \Longleftrightarrow(1,1) \text { is odd } \\
& \Longleftrightarrow 2(1,1) \equiv 2(\bmod 4) \Longleftrightarrow M+L \equiv 0(\bmod 4)
\end{aligned}
$$

To prove the final assertion of the theorem, notice that if 2 is a cube mod $p$, so that $L$ and $M$ are even, then $p=L^{\prime 2}+27 M^{\prime 2}$ with $L^{\prime}:=L / 2$ and $M^{\prime}:=M / 2$. Conversely, if $p=L^{\prime 2}+27 M^{\prime 2}$ for some integers $L^{\prime}$ and $M^{\prime}$, then $4 p=L^{2}+27 M^{2}$ where $L=2 L^{\prime}$ and $M=2 M^{\prime}$. Since the integers $L$ and $M$ in such a representation are uniquely determined up to sign, it follows that $L$ and $M$ are even in all such representations, so that 2 is a cube modulo $p$.

Theorem 2.27 (Cubic character of 3). Under the same assumptions as the previous theorem,

3 is a cube modulo $p \Longleftrightarrow 3 \mid M$,

$$
\begin{aligned}
\operatorname{ind}_{g}(3) \equiv 1(\bmod 3) & \Longleftrightarrow M \equiv-1(\bmod 3) \\
\operatorname{ind}_{g}(2) & \equiv 2(\bmod 3) \Longleftrightarrow M \equiv+1(\bmod 3)
\end{aligned}
$$

Proof. As $\beta$ runs through all the elements of $\mathbf{F}_{p}^{\times} \backslash\{1\}$, the expression $(\beta-1)^{-1}$ assumes all the values in $\mathbf{F}_{p}^{\times} \backslash\{-1\}$. So by Wilson's theorem, $\prod_{\beta \in \mathbf{F}_{p}^{\times} \backslash\{1\}}(\beta-1)^{-1}=1$. With $g$ our chosen generator, we put $\omega:=g^{(p-1) / 3}$ and let $H:=\left\{1, \omega, \omega^{2}\right\}$ be the subgroup of $\mathbf{F}_{p}^{\times}$generated by $\omega$. Let $A:=$ $\left\{\gamma_{1}=1, \gamma_{2}, \ldots, \gamma_{(p-1) / 3}\right\}$ be a complete set of coset representatives for $H$. Then we have

$$
\begin{aligned}
\prod_{\substack{\beta \in(\mathbf{Z} / p \mathbf{Z})^{\times} \\
\beta \neq 1}} \frac{1}{\beta-1} & =\frac{1}{\omega-1} \frac{1}{\omega^{2}-1} \prod_{1 \neq \gamma \in A} \frac{1}{\gamma-1} \frac{1}{\gamma \omega-1} \frac{1}{\gamma \omega^{2}-1} \\
& =\frac{1}{1-\omega} \frac{1}{1-\omega^{2}} \prod_{1 \neq \gamma \in A} \frac{1}{\gamma-1} \frac{1}{\gamma-\omega} \frac{1}{\gamma-\omega^{2}} \\
& =\frac{1}{3} \prod_{1 \neq \gamma \in A} \frac{1}{\gamma^{3}-1} .
\end{aligned}
$$

As $\gamma$ runs through the elements of $A \backslash\{1\}$, the element $\gamma^{3}-1$ runs exactly once through the immediate predecessors of every cube $\neq 1$. It follows that

$$
0=\operatorname{ind}_{g}(1) \equiv-\operatorname{ind}_{g}(3)-\sum_{1 \neq \gamma \in A} \operatorname{ind}_{g}\left(\gamma^{3}-1\right)(\bmod p-1) ;
$$

modulo 3 this implies that

$$
-\operatorname{ind}_{g}(3)-0(0,0)-1(1,0)-2(2,0) \equiv 0(\bmod 3),
$$

i.e.,

$$
\operatorname{ind}_{g}(3) \equiv-(1,0)-2(2,0)=-b-2 c \equiv c-b \equiv-M(\bmod 3),
$$

as we sought to show.
Remark. Gauss's first proof of Theorem 2.27 (which has been preserved in [Gau73a, pp. 10-11]) was a good deal more intricate. The elegant argument described above was discovered subsequently by Gauss, and recorded on January 6th, 1809 in his mathematical diary:

The theorem for the cubic residue 3 is proved with an elegant special method by considering the values of $\frac{x+1}{x}$ where three each always have the values $a, a \epsilon, a \epsilon^{2}$, with the exception of two which give $\epsilon, \epsilon^{2}$, but these are

$$
\frac{1}{\epsilon-1}=\frac{\epsilon^{2}-1}{3}, \quad \frac{1}{\epsilon^{2}-1}=\frac{\epsilon-1}{3}
$$

with product $\equiv \frac{1}{3}$.
For many years this comment remained obscure. The reconstruction presented here is due to Gröger [Grö06].
7.3. Jacobi's rational cubic reciprocity law. We now show how to derive Jacobi's original form of cubic reciprocity from Kummer's criterion (Theorem 2.15) and our determination of the cubic period polynomial. Sun's version of Jacobi's law is treated in §7.4.

First we recall the statement of Jacobi's law:
Theorem 2.28 (Jacobi). Let $p$ and $q$ be distinct primes with $p, q>3$ and $p \equiv 1(\bmod 3)$. Write $4 p=L^{2}+27 M^{2}$. Then

$$
\begin{equation*}
q \text { is a cube in } \mathbf{F}_{p} \Longleftrightarrow \frac{L+3 M \sqrt{-3}}{L-3 M \sqrt{-3}} \text { is a cube in } \mathbf{F}_{q}(\sqrt{-3}) \text {. } \tag{2.20}
\end{equation*}
$$

We can (and do) assume for the proof of Theorem 2.28 that the sign of $L$ is chosen so that $L \equiv 1(\bmod 3)$. Indeed, replacing $L$ with $-L$ has the effect of replacing the ratio on the right-hand side of (2.20) with its reciprocal, and this new ratio is a cube exactly when the original is.

Let $\phi$ (respectively $\hat{\phi}$ ) be the cubic period polynomial (respectively reduced period polynomial) whose coefficients were determined in §7.1. Then

$$
\operatorname{Disc}(\hat{\phi})=4(3 p)^{3}-27(p L)^{2}=27 p^{2}\left(4 p-L^{2}\right)=3^{6} p^{2} M^{2}
$$

$\operatorname{But} \operatorname{Disc}(\hat{\phi})=3^{6} \cdot \operatorname{Disc}(\phi)$, so that

$$
\operatorname{Disc}(\phi)=p^{2} M^{2} .
$$

Since $e=3$ is prime, part (iii) of Kummer's criterion (Theorem 2.15) yields the following special case of Theorem 2.28. (Note that if $q \mid M$, then the quotient on the right-hand side of $(2.20)$ is $L / L=1$, which is a cube in $\mathbf{F}_{q}(\sqrt{-3})$.)

Lemma 2.29. Let $p$ and $q$ be distinct primes with $p, q>3$ and $p \equiv$ $1(\bmod 3)$. Write $4 p=L^{2}+27 M^{2}$ with $L \equiv 1(\bmod 3)$. If $q$ divides $M$, then $q$ is a cube in $\mathbf{F}_{p}$.

It remains to treat the case when $q>3$ and $q \nmid p M$. Here we use Corollary 2.16:
$q$ is a cube modulo $p \Longleftrightarrow \phi$ has a root $\bmod q \Longleftrightarrow \hat{\phi}$ has a root $\bmod q$,
the last implication following from (2.6). To analyze when $\hat{\phi}$ has a root in $\mathbf{F}_{q}$, we use the classical solution of the cubic equation.
$\star$ Theorem 2.30 (Cardano). Let $f(T)=T^{3}+a T-b$ be a cubic polynomial with coefficients in a field $F$ of characteristic $\neq 2,3$. Suppose also that $a \neq 0$. Then the roots of $f$ in an algebraic closure of $F$ are given by

$$
w+\frac{-a / 3}{w}, \quad \text { where } \quad w^{3}=\frac{b}{2} \pm \sqrt{\frac{b^{2}}{4}+\frac{a^{3}}{27}},
$$

where $w$ ranges over all six cube roots corresponding to the two choices of sign.

Applied to our situation we find:
Corollary 2.31. Let $p \equiv 1(\bmod 3)$ and let $q>3$ be a prime not dividing $p M$. Then the roots of the reduced cubic period polynomial $T^{3}-3 p T-p L$ in an algebraic closure of $\mathbf{F}_{q}$ can be described by

$$
w+\frac{p}{w}, \quad \text { where } \quad w^{3}=p \frac{L \pm 3 M \sqrt{-3}}{2} .
$$

Let $w$ be one of these cube roots. Since the elements of $\mathbf{F}_{q}$ can be characterized as the fixed points of the $q$ th power map, for the root corresponding to $w$ we have

$$
\begin{aligned}
w+p / w \in \mathbf{F}_{q} & \Longleftrightarrow(w+p / w)^{q}=(w+p / w) \\
& \Longleftrightarrow w^{q}+p / w^{q}=w+p / w .
\end{aligned}
$$

To analyze the last of these equivalent statements, we use the following lemma, whose proof is left as Exercise 3.

Lemma 2.32. Let $F$ be a field of characteristic other than $p$. If $x, y \in F$ and $x+p / x=y+p / y$, then either $x=y$ or $x=p / y$.

We conclude that if $w+p / w \in \mathbf{F}_{q}$, then either $w^{q}=w$ or $w^{q}=p / w$. We now show that the first possibility can only occur if $q \equiv 1(\bmod 3)$ and that the latter can only occur if $q \equiv 2(\bmod 3)$.

Lemma 2.33. Let $p$ and $q$ be distinct primes with $p, q>3$ and $p \equiv$ $1(\bmod 3)$. Suppose $q \nmid p M$. Suppose the element $w$ in a fixed algebraic closure of $\mathbf{F}_{q}$ satisfies

$$
\begin{equation*}
w^{3}=p \frac{L \pm 3 M \sqrt{-3}}{2} \in \mathbf{F}_{q}(\sqrt{-3}) \tag{2.21}
\end{equation*}
$$

for some choice of sign. Then

$$
\begin{aligned}
& w^{3 q}=w^{3} \text { if and only if } q \equiv 1(\bmod 3) \\
& \qquad \text { while } \quad w^{3 q}=p^{3} / w^{3} \text { if and only if } q \equiv 2(\bmod 3) .
\end{aligned}
$$

Consequently, $w^{q}=w$ implies $q \equiv 1(\bmod 3)$ and $w^{q}=p / w$ implies $q \equiv$ $2(\bmod 3)$.

Proof. We have

$$
\begin{equation*}
w^{3 q}=\left(w^{3}\right)^{q}=p^{q}\left(\frac{L \pm 3 M \sqrt{-3}}{2}\right)^{q}=p \frac{L \pm 3 M\left(\frac{-3}{q}\right) \sqrt{-3}}{2} . \tag{2.22}
\end{equation*}
$$

As $M \neq 0$ in $\mathbf{F}_{q}$ by hypothesis, the right-hand side agrees with $w^{3}$ exactly when $\left(\frac{-3}{q}\right)=1$, i.e., when $q \equiv 1(\bmod 3)$. Since

$$
p^{3} / w^{3}=\frac{p^{3}}{p(L \pm 3 M \sqrt{-3}) / 2}=p \frac{p}{(L \pm 3 M \sqrt{-3}) / 2}=p \frac{L \mp 3 M \sqrt{-3}}{2},
$$

the right-hand side of $(2.22)$ agrees with $p^{3} / w^{3}$ exactly when $\left(\frac{-3}{q}\right)=-1$, i.e., when $q \equiv 2(\bmod 3)$.

We prove Theorem 2.28 by analyzing for which primes $p \equiv 1(\bmod 3)$ we have $w^{q}=w$ and for which primes $p \equiv 2(\bmod 3)$ we have $w^{q}=p / w$. By Lemma 2.29, we can assume in these proofs that $q \nmid M$.

In what follows we let $\sqrt{-3}$ denote a fixed square root of -3 in an algebraic closure of $\mathbf{F}_{q}$ and we let $w$ be an element of this algebraic closure satisfying (2.21). For notational convenience we also set

$$
\pi:=\frac{L \pm 3 M \sqrt{-3}}{2} \quad \text { and } \quad \pi^{\prime}:=\frac{L \mp 3 M \sqrt{-3}}{2}
$$

so that $\pi \pi^{\prime}=p$ and $w^{3}=p \pi$.

Proof of the Jacobi law for $q \equiv 1(\bmod 3)$. In this case

$$
\begin{aligned}
w+p / w \in \mathbf{F}_{q} & \Longleftrightarrow w^{q}=w \Longleftrightarrow w^{q-1}=1 \\
& \Longleftrightarrow(p \pi)^{(q-1) / 3}=1 \Longleftrightarrow\left(\pi^{2} \pi^{\prime}\right)^{(q-1) / 3}=1
\end{aligned}
$$

Since $q \equiv 1(\bmod 3)$, we have $\mathbf{F}_{q}(\sqrt{-3})=\mathbf{F}_{q}$. Hence $\pi$ and $\pi^{\prime}$ are elements of $\mathbf{F}_{q}$ (and are nonzero since they multiply to the nonzero element $p$ ). So by Euler's criterion, the above holds

$$
\Longleftrightarrow \pi^{2} \pi^{\prime} \text { is a cube in } \mathbf{F}_{q} \Longleftrightarrow \frac{\pi^{2} \pi^{\prime}}{\pi^{\prime 3}}=\pi^{\prime} / \pi \text { is a cube in } \mathbf{F}_{q}
$$

If the minus sign holds in the definition of $\pi$, then this is exactly the criterion appearing in (2.20). If the plus sign holds, then we have only to note that $\pi / \pi^{\prime}$ is a cube if and only if $\pi^{\prime} / \pi$ is a cube, and we again recover Jacobi's criterion.

Since this computation was valid for any choice of $w$, we have proved more than required: We have shown that if the right-hand side of (2.20) is a cube in $\mathbf{F}_{q}(\sqrt{-3})$, then the reduced period polynomial has all its roots (not just one) defined modulo $q$. Conversely, if this quotient is not a cube, then none of the roots of the reduced period polynomial lie in $\mathbf{F}_{q}$.

Proof of the Jacobi law for $q \equiv 2(\bmod 3)$. In this case

$$
w+p / w \in \mathbf{F}_{q} \Longleftrightarrow w^{q}=p / w \Longleftrightarrow w^{q+1}=p .
$$

By Lemma 2.33, we have $w^{3(q+1)}=p^{3}$. Since the cube roots of unity lie outside $\mathbf{F}_{q}$,

$$
w^{q+1}=p \Leftrightarrow w^{q+1} \in \mathbf{F}_{q} \Leftrightarrow w^{(q+1)(q-1)}=1 \Leftrightarrow p^{\left(q^{2}-1\right) / 3} \pi^{\left(q^{2}-1\right) / 3}=1 .
$$

But for a nonzero $\alpha \in \mathbf{F}_{q}(\sqrt{-3})$, we have $\alpha^{\left(q^{2}-1\right) / 3}=1$ precisely when $\alpha$ is a cube. Note that since $q \equiv 2(\bmod 3)$, every element of $\mathbf{F}_{q}$ (in particular, the element $p$ ) is a cube in both $\mathbf{F}_{q}$ and $\mathbf{F}_{q}(\sqrt{-3})$. Hence

$$
\begin{aligned}
p^{\left(q^{2}-1\right) / 3} \pi^{\left(q^{2}-1\right) / 3}=1 & \Longleftrightarrow \pi^{\left(q^{2}-1\right) / 3}=1 \\
& \Longleftrightarrow \pi \text { is a cube in } \mathbf{F}_{q}(\sqrt{-3}) \\
& \Longleftrightarrow \pi^{2} \text { is a cube in } \mathbf{F}_{q}(\sqrt{-3}) \\
& \Longleftrightarrow \pi^{2} / p=\pi / \pi^{\prime} \text { is a cube in } \mathbf{F}_{q}(\sqrt{-3}) .
\end{aligned}
$$

The proof is now completed as in the case $q \equiv 1(\bmod 3)$.
7.4. Sun's form of Jacobi's law. We now prove Sun's pretty equivalent form of Jacobi's law (see [Sun98]), enunciated as the second half of Theorem 2.5 in the introduction. Recall that for each prime $q>3$ we defined the group $G=G(q)$ by

$$
G=\left\{[a, b]: a, b \in \mathbf{F}_{q}, a^{2}+3 b^{2} \neq 0\right\},
$$

where we identify $[a, b]$ and $[c, d]$ if $a=\lambda c, b=\lambda d$ for some nonzero $\lambda \in \mathbf{F}_{q}$, and where we multiply according to the rule

$$
[a, b] \odot[c, d]=[a c-3 b d, a d+b c] .
$$

All of the group axioms are quickly verified, with $[1,0]$ as the identity element, except associativity. We leave this to the reader to check by a direct calculation.
Lemma 2.34. We have $\# G=q-\left(\frac{-3}{q}\right)$.
Proof. Every element besides $[1,0]$ can be written uniquely in the form $[a, 1]$ with $a \in \mathbf{F}_{q}$. We have $[a, 1] \in G$ if and only if $a^{2} \neq-3$. Hence

$$
\begin{aligned}
\# G & =1+\# \mathbf{F}_{q}-\#\left\{a \in \mathbf{F}_{q}: a^{2}=-3\right\} \\
& =1+q-\left(1+\left(\frac{-3}{q}\right)\right)=q-\left(\frac{-3}{q}\right) .
\end{aligned}
$$

Lemma 2.35. Let $\psi$ be the map from $G$ to $\mathbf{F}_{q}(\sqrt{-3})^{\times}$defined by

$$
\psi([a, b]):=\frac{a+b \sqrt{-3}}{a-b \sqrt{-3}} .
$$

Then $\psi$ is an injective homomorphism. Hence $G$ is cyclic.
Proof. We need to check first that $\psi$ is well-defined: This follows because $a^{2}+3 b^{2} \neq 0$ and because we are taking a ratio on the right-hand side (so that the ambiguity in $[a, b]$ up to scaling disappears). To see that $\psi$ is a homomorphism, we compute:

$$
\begin{aligned}
\psi([a, b] \odot[c, d]) & =\psi([a c-3 b d, a d+b c]) \\
& =\frac{a c-3 b d+(a d+b c) \sqrt{-3}}{a c-3 b d-(a d+b c) \sqrt{-3}} \\
& =\frac{a+b \sqrt{-3}}{a-b \sqrt{-3}} \cdot \frac{c+d \sqrt{-3}}{c-d \sqrt{-3}}=\psi([a, b]) \psi([c, d]) .
\end{aligned}
$$

To see that $\psi$ is injective, it suffices to prove that its kernel is trivial: But

$$
\psi([a, b])=1 \Longrightarrow \frac{a+b \sqrt{-3}}{a-b \sqrt{-3}}=1
$$

and this implies that $b=0$. Hence $[a, b]=[1,0]$ is the identity of $G$. This proves $\psi$ is an embedding as claimed.

The cyclicity of $G$ is an easy corollary: We can view $G$ as a subgroup of $\mathbf{F}_{q}(\sqrt{-3})^{\times}$, and every finite subgroup of the multiplicative group of a field is cyclic.

We can now prove Sun's form of Jacobi's law:

Theorem 2.36. Let $p$ and $q$ be distinct primes, with $p, q>3$ and $p \equiv$ $1(\bmod 3)$. Write $4 p=L^{2}+27 M^{2}$ with integers $L$ and $M$, and let $G=G(q)$ be the group defined above. Then

$$
q \text { is a cube modulo } p \Longleftrightarrow[L, 3 M] \text { is a cube in } G \text {. }
$$

Proof. Let $H$ be the image of $\psi$, where $\psi$ is the map of Lemma 2.35 (so that $\# H=\# G)$. By Theorem 2.28,

$$
\begin{aligned}
q \text { is a cube modulo } p & \Longleftrightarrow \psi([L, 3 M]) \text { is a cube in } \mathbf{F}_{q}(\sqrt{-3}) \\
& \Longleftrightarrow \psi([L, 3 M])^{\# \mathbf{F}_{q}(\sqrt{-3})^{\times} / 3}=1 \\
& \Longleftrightarrow \psi([L, 3 M])^{\operatorname{gcd}\left(\# H, \# \mathbf{F}_{q}(\sqrt{-3})^{\times} / 3\right)}=1 \\
& \Longleftrightarrow \psi([L, 3 M])^{\# H / 3}=1 \\
& \Longleftrightarrow \psi\left([L, 3 M]^{\# H / 3}\right)=1 .
\end{aligned}
$$

Since $\psi$ has trivial kernel, the last equality holds precisely when $[L, 3 M]^{\# H / 3}$ is the identity of $G$. Since $\# H=\# G$, this holds if and only if $[L, 3 M]$ is a cube in $G$.

As we mentioned in the introduction, Jacobi's cubic reciprocity law implies that whether $q$ is a residue or nonresidue of $p$ depends only on the ratio $L / M \bmod q$. These ratios are the subject of the following two theorems. We leave their proofs as Exercises 16 and 17.

Theorem 2.37 (Cunningham \& Gosset [CG20]). Let $p \equiv 1(\bmod 3)$ be prime and write $4 p=L^{2}+27 M^{2}$ with integers $L$ and $M$. Let $q>3$ be a prime distinct from $p$, and let $n=\frac{1}{3}\left(q-\left(\frac{-3}{q}\right)\right)$. Then $q$ is a cube $\bmod p$ if and only if

$$
\sum_{\substack{0 \leq j \leq n \\ j \equiv 1(\bmod 2)}} 3^{j}(-3)^{(j-1) / 2}\binom{n}{j} L^{n-j} M^{j} \equiv 0(\bmod q) .
$$

A more explicit description of these ratios is provided by the next result:
Theorem $2.38(\operatorname{Sun})$. Let $p \equiv 1(\bmod 3)$ be prime and write $4 p=L^{2}+$ $27 M^{2}$ with integers $L$ and $M$. Let $q>3$ be a prime distinct from $p$. Then $q$ is a cubic residue modulo $p$ if and only if either $q$ divides $M$ or $\frac{L}{3 M} \equiv$ $\frac{x^{3}-9 x}{3 x^{2}-3}(\bmod q)$ for some integer $x$.

## Notes

Jacobi's law (Theorem 2.5) is an example of a rational reciprocity law; the word "rational" is here because the statement of the law refers only to rational integers. This is in contrast to Eisenstein's cubic reciprocity law,
which is not a statement about rational primes but a statement about primes in the ring $\mathbf{Z}[\omega]$, where $\omega$ is a complex primitive cube root of unity.

While Eisenstein's law is harder to state, it has the advantage of being applicable to more problems. To see why Jacobi's law is not the end of the story (even if one is concerned just with $\mathbf{Z}$ and not $\mathbf{Z}[\omega]$ ), consider the problem of determining the primes $p \equiv 1(\bmod 3)$ for which 35 is a cube modulo $p$. Theorem 2.5 suffices to tell us when 5 is a cube modulo $p$ and when 7 is a cube modulo $p$. But if neither 5 nor 7 are cubes modulo $p$, the status of 35 is still undetermined: In this case whether or not 35 is a cube modulo $p$ depends on whether 5 and 7 belong to the same coset or different cosets of $\left(\mathbf{F}_{p}^{\times}\right)^{3}$ in $\mathbf{F}_{p}^{\times}$.

This suggests the following: Given a prime $q$ different than $p$, we would like to know not merely when $q^{(p-1) / 3}=1$, but which cube root of unity $q^{(p-1) / 3}$ represents in $\mathbf{F}_{p}$; this is not a question that Jacobi's law answers. However, an answer can be coaxed out of Eisenstein's law. This requires one to translate the problem into the setting of $\mathbf{Z}[\omega]$, where Eisenstein's law operates, work out the answer, and then translate back! Luckily, the heavy lifting has been done by Sun ([Sun98, Corollary 2.1, Theorem 2.2]; see also the paper of von Lienen [vL79]). He proves the following:
$\star$ Theorem 2.39. Let $p, q>3$ be distinct primes, and suppose $p \equiv$ $1(\bmod 3)$. Write $4 p=L^{2}+27 M^{2}$. Put

$$
\omega:=\frac{-1-L / 3 M}{2} ;
$$

by the choice of $L$ and $M$, this represents a primitive cube root of unity in $\mathbf{F}_{p}$. Write $\bar{\omega}$ for the element $[1,1]$ of $G(q)$, where $G(q)$ is the group considered in Sun's Theorem 2.36; note that $\bar{\omega}$ is an element of order 3 in $G(q)$. For each $i \in\{0,1,2\}$, we have

$$
q^{\frac{p-1}{3}} \equiv \omega^{i}(\bmod p) \Longleftrightarrow[L, 3 M]^{\frac{q-\left(\frac{-3}{q}\right)}{3}}=\bar{\omega}^{i} \quad \text { in } G(q) .
$$

Note that when $i=0$, this reduces to Theorem 2.36. For an excellent account of Eisenstein's cubic reciprocity law, see Chapter 9 of the text of Ireland \& Rosen [IR90] or Chapter 7 of Lemmermeyer's beautiful monograph [Lem00]. For further discussion of rational reciprocity laws, see [Lem00, Chapter 5] and [BEW98, Chapters 7 and 8].

## Exercises

1. Let $p \equiv 1(\bmod 3)$. Suppose that $A_{1}, B_{1}, A_{2}, B_{2} \in \mathbf{Z}$ and $A_{1}^{2}+27 B_{1}^{2}=$ $A_{2}^{2}+27 B_{2}^{2}=4 p$. Prove that $A_{1}= \pm A_{2}$ and $B_{1}= \pm B_{2}$. Hint: Verify the identity

$$
16 p^{2}=\left(A_{1} A_{2} \pm 27 B_{1} B_{2}\right)^{2}+27\left(A_{2} B_{1} \mp A_{1} B_{2}\right)^{2} .
$$

Also, check that

$$
p \mid\left(A_{1} A_{2}-27 B_{1} B_{2}\right)\left(A_{2} B_{1}-A_{1} B_{2}\right)
$$

and

$$
p \mid\left(A_{1} A_{2}+27 B_{1} B_{2}\right)\left(A_{2} B_{1}+A_{1} B_{2}\right) .
$$

Deduce that $p \mid A_{2} B_{1} \pm A_{1} B_{2}$ for one of the choices of sign, and conclude that $A_{1} / A_{2}= \pm B_{1} / B_{2}$.
2. Show that if $\varphi(n)$ is a power of 2 , then $n$ has the form $2^{e} P$, where $e \geq 0$ and $P$ is a product (possibly empty) of distinct Fermat primes.
3. Prove Lemma 2.32.
4. Say that $\alpha \in \mathbf{R}$ is real-constructible if it is possible to construct two points a distance $|\alpha|$ apart.
(a) Prove (or look up) the following (geometric) lemma: If $\alpha$ and $\beta$ are two real-constructible numbers, then so are

$$
\alpha \pm \beta, \quad \alpha \beta, \quad 1 / \alpha \quad(\text { if } \alpha \neq 0), \quad \sqrt{\alpha} \quad(\text { if } \alpha \geq 0) .
$$

Hence the real-constructible numbers form a subfield of $\mathbf{R}$, say Cons $_{\mathbf{R}}$. Show, moreover, that the point $(x, y)$ is constructible if and only if its components $x$ and $y$ are both real-constructible.
(b) Suppose we have a tower of subfields of the real numbers

$$
\mathbf{Q}:=K_{0} \subset K_{1} \subset K_{2} \subset \cdots \subset K_{m}
$$

where $\alpha \in K_{m}$ and, for $i>0$, each $K_{i}=K_{i-1}\left(\sqrt{\beta_{i}}\right)$ for some nonnegative $\beta_{i} \in K_{i-1}$. Using part (a), prove that $\alpha$ is realconstructible.
(c) Let $L$ be the line described by the equation $a x+b y=c$, and let $C$ be the circle described by the equation $\left(x-x_{0}\right)^{2}+\left(y-y_{0}\right)^{2}=r^{2}$. Let $K=\mathbf{Q}\left(a, b, c, x_{0}, y_{0}, r\right)$. Prove that each coordinate of a point of intersection of $L$ and $C$ lies either in $K$ or in a quadratic extension of $K$.
(d) Use (c) to prove the converse of (b): If $\alpha$ is real-constructible, then there is such a tower whose last term contains $\alpha$.
(e) Prove that the point $(x, y)$ is constructible if and only if $x+i y \in$ $\operatorname{Cons}_{\mathbf{R}}(i)$. Now prove that the elements of (the field!) $\mathrm{Cons}_{\mathbf{R}}(i)$ are exactly the elements described in Lemma 2.6. For one containment you may find helpful the identity

$$
\sqrt{x+i y}=\frac{1}{2} \sqrt{2}\left(\sqrt{\sqrt{x^{2}+y^{2}}+x}+i \operatorname{sgn}(y) \sqrt{\sqrt{x^{2}+y^{2}}-x}\right) .
$$

Here $\operatorname{sgn}(y) \in\{0,1,-1\}$ is defined as $y /|y|$ for $y \neq 0$ and defined to be 0 when $y=0$.
5 . Prove that the following are equivalent for every $n \geq 3$ :
(a) It is possible to construct all the vertices of a regular $n$-gon,
(b) Some primitive $n$th root of unity is constructible,
(c) Every primitive $n$th root of unity is constructible.
6. (Gauss [Gau86, Art. 354]) In this exercise we make explicit Theorem 2.8 for the case $p=17$. We use the notation of Figure 1 for the Gaussian periods.
(a) Using Lemma 2.19, prove the polynomial identities
(i) $(T-(8,1))(T-(8,3))=T^{2}+T-4$,
(ii) $(T-(4,1))(T-(4,9))=T^{2}-(8,1) T-1$,
(iii) $(T-(4,3))(T-(4,10))=T^{2}-(8,3) T-1$,
(iv) $(T-(2,1))(T-(2,13))=T^{2}-(4,1) T+(4,3)$,
(v) $(T-(1,1))(T-(1,16))=T^{2}-(2,1) T+1$.
(b) Show that one can choose the primitive 17 th root of unity $\zeta$ so that

$$
(8,1)=\frac{-1+\sqrt{17}}{2} \quad \text { and } \quad(4,1)=\frac{(8,1)+\sqrt{(8,1)^{2}+4}}{2}
$$

Of course the difficulty is in proving that we can make the plus sign hold in both places.
(c) The choices of sign in (b) force a choice of sign for $(4,3)$ : To see this, prove that

$$
((4,1)-(4,9))((4,3)-(4,10))=2((8,1)-(8,3))>0,
$$

and deduce that $(4,3)=\frac{1}{2}\left((8,3)+\sqrt{(8,3)^{2}+4}\right)$.
(d) Prove that we can choose $\zeta$ as in (b) so that

$$
(2,1)=\frac{(4,1)+\sqrt{(4,1)^{2}-4(4,3)}}{2}
$$

again, the nontrivial aspect is to prove that we can force the plus sign. (Note that $(4,1)^{2}-4(4,3)>0$, as follows from a rough numerical calculation.)
(e) We have

$$
(2,1)=\zeta+\zeta^{g^{8}}=\zeta+\zeta^{-1}=2 \Re(\zeta)
$$



Figure 2. Diagram accompanying Richmond's construction of the 17gon (see Exercise 7), based on [HW08, Fig. 5, p. 76].

Obtain a rough numerical approximation (on a calculator, say) of $(2,1)$ sufficient to prove to pin down $\zeta$ to one of the two values $e^{ \pm 2 \pi i / 17}$; hence $(2,1)=2 \cos \frac{2 \pi}{17}$.
(f) Prove that $e^{2 \pi i / 17}$ and $e^{-2 \pi i / 17}$ are the roots of $T^{2}-(2,1) T+1$.
(g) Combining (a)-(e), show that
$(2,1)=2 \cos \frac{2 \pi}{17}=\frac{1}{8} \sqrt{34-2 \sqrt{17}}-\frac{1}{8}+\frac{1}{8} \sqrt{17}+$
$\frac{1}{8} \sqrt{68+12 \sqrt{17}-2 \sqrt{34-2 \sqrt{17}}+2 \sqrt{34-2 \sqrt{17}} \sqrt{17}-16 \sqrt{34+2 \sqrt{17}}}$.
Now use (f) to compute an explicit representation of $\zeta_{17}$. (You may wish to use a computer algebra system for this part.)
Lecture 7 of [Rad64] is a self-contained account of the results of this exercise; see also Hardy \& Wright [HW08, §5.8].
7. The result of the preceding exercise gives us an explicit way of constructing the 17 -gon; however, such a direct attack is both inefficient and onerous. In 1893, Richmond proposed the following alternative geometric construction ([Ric93, Ric09]):

Let $O A, O B$ [Figure 2] be two perpendicular radii of a circle. Make $O I$ one-fourth of $O B$, and the angle $O I E$
one-fourth of $O I A$; also find in $O A$ produced a point $F$ such that $E I F$ is $45^{\circ}$. Let the circle on $A F$ as diameter cut $O B$ in $K$, and let the circle whose centre is $E$ and radius $E K$ cut $O A$ in $N_{3}$ and $N_{5}$; then if ordinates $N_{3} P_{3}$, $N_{5} P_{5}$ are drawn to the circle, the arcs $A P_{3}, A P_{5}$ will be $3 / 17$ and $5 / 17$ of the circumference.
Prove Richmond's assertions. If you have trouble with this, Hardy \& Wright [HW08, §5.8] present his construction in detail.
8. (Luca [Luc00b]) Say that the natural number $n \geq 2$ has property (C) if both the regular ( $n-1$ )-gon and regular $n$-gon are constructible. Using the Gauss-Wantzel Theorem, show that if $n$ has property (C), then either $n$ is a Fermat prime or $n \in\left\{2 \cdot 3,2^{2}, 2^{2^{2}}, 2^{2^{3}}, 2^{2^{4}}, 2^{2^{5}}\right\}$. Proceed as follows:
(a) Consider a nonempty product of distinct Fermat numbers $F_{m}=$ $2^{2^{m}}+1$, say

$$
\begin{equation*}
F_{n_{0}} F_{n_{1}} \cdots F_{n_{k-1}}, \tag{2.23}
\end{equation*}
$$

where $0 \leq n_{0}<n_{1}<\cdots<n_{k-1}$.
(i) Prove that this product has precisely $2^{k}$ nonzero digits in its binary expansion.
(ii) Show that, moreover, there are $1+2^{n_{0}}+2^{n_{1}}+\cdots+2^{n_{k-1}}$ total binary digits in this product. Thus, if we start with the number of binary digits in the product, subtract one and compute the binary expansion, we can read off the $n_{i}$ corresponding to the Fermat number factors.
(b) Using (a), prove that any odd number $n$ with property (C) is a Fermat prime.
(c) Suppose $n$ is even and has property (C). Using (b), show that if $n \equiv 2(\bmod 4)$, then $n-1=F_{1}$, and so $n=6$.
(d) Finally, suppose $n$ has property $C$ where $4 \mid n$. Since $n-1 \equiv$ $3(\bmod 4)$, if we write $n-1$ in the form (2.23), then $n_{0}=0$. Suppose that $n_{0}=0, n_{1}=1, n_{2}=2, \ldots, n_{k^{\prime}}=k^{\prime}$ for a certain $k^{\prime} \geq 0$ while $n_{j} \geq k^{\prime}+2$ for the remaining indices $k^{\prime}<j<k$. Verify that in this case the binary expansion of $n-1$ ends with precisely $2^{k^{\prime}+1}$ trailing 1 's, and the binary expansion of $n$ contains precisely $2^{k}-2^{k^{\prime}+1}+1$ nonzero binary digits.
Now obtain a contradiction to (i) unless $k=k^{\prime}+1$, i.e., unless $n-1=F_{0} F_{1} \cdots F_{k^{\prime}}$. Complete the proof making use of Euler's discovery that $F_{5}$ is composite.
9. Here we give two proofs for the irreducibility of the cyclotomic polynomials $\Phi_{n}(T)$ : Let $\zeta$ be a primitive $n$th root of unity and let $f(T) \in \mathbf{Q}[T]$ be its minimal polynomial. It is easy to show that $f(T) \in \mathbf{Z}[T]$ and that
$f(T)$ divides $\Phi_{n}(T)$ in $\mathbf{Z}[T]$. We would like to show that $f(T)=\Phi_{n}(T)$, and for this it suffices to prove that $\zeta^{a}$ is a root of $f$ for each $a$ coprime to $n$.
(a) Prove that a nonzero element of the ring $\mathbf{Z}[\zeta]$ is divisible by only finitely many rational primes $p$.
(b) Prove that $p \mid f\left(\zeta^{p}\right)$ in $\mathbf{Z}[\zeta]$ for every prime $p$ not dividing $n$.
(c) (Grandjot [Gra23]) We can now give a simple proof by means of Dirichlet's theorem. Let $a$ be coprime to $n$. Letting $p$ run through the primes congruent to $a \bmod n$, show that the single element $f\left(\zeta^{a}\right)$ has infinitely many rational prime divisors; conclude from part (i) that $f\left(\zeta^{a}\right)=0$ as desired.
(d) (Landau [Lan28]) Here is an alternative argument avoiding Dirichlet's result. Using (a), show that we can choose a number $B$ (depending only on $n$ ) so that if $p>B$ is prime and $a$ is coprime to $n$, then either $f\left(\zeta^{a}\right)=0$ or $p \nmid f\left(\zeta^{a}\right)$.
Fix such a $B$, and fix a particular integer $a$ coprime to $n$. Choose a positive integer $m$ with $m \equiv a(\bmod n)$ and $m$ coprime to $\prod_{p \leq B} p$. Factor $m=q_{1} q_{2} \cdots q_{j}$ as a product of primes, and show successively that all of

$$
\zeta^{q_{1}}, \zeta^{q_{1} q_{2}}, \ldots, \zeta^{q_{1} \cdots q_{j}}=\zeta^{a}
$$

are roots of $f$.
10. (Ankeny; see [Ank60]) Fix a prime $e$. Let $p$ and $q$ be primes distinct from each other and distinct from $e$ with $p \equiv 1(\bmod e)$. Let $\zeta_{e}$ and $\zeta_{p}$ be fixed primitive $e$ th and $p$ th roots of unity in a fixed algebraic closure $\overline{\mathbf{F}}_{q}$ of $\mathbf{F}_{q}$. Let $\chi: \mathbf{F}_{p}^{\times} \rightarrow \overline{\mathbf{F}}_{q}^{\times}$be a homomorphism whose image is precisely the set of $e$ th roots of unity in $\overline{\mathbf{F}}_{q}^{\times}$. We define the Gauss sum $\tau_{a}(\chi)$ by

$$
\tau_{a}(\chi):=\sum_{n=1}^{p-1} \chi(n) \zeta_{p}^{a n}
$$

If $a=1$, we write $\tau_{1}(\chi)=\tau(\chi)$.
(a) Prove that $\tau_{a}(\chi) \tau_{-a}\left(\chi^{-1}\right)=p$ for every $a$ not divisible by $p$. So, in particular, $\tau_{a}(\chi)$ is nonzero for all such $a$. Hint:

$$
\tau_{a}(\chi) \tau_{-a}\left(\chi^{-1}\right)=\sum_{n, m \in \mathbf{F}_{p}^{\times}} \chi\left(n m^{-1}\right) \zeta_{p}^{a(n-m)}=\sum_{l \in \mathbf{F}_{p}^{\times}} \chi(l) \sum_{m \in \mathbf{F}_{p}^{\times}} \zeta_{p}^{a m(l-1)} .
$$

(b) Let $f$ be the order of $q(\bmod e)$. Prove that $\tau(\chi)^{q^{f}}=\chi(q)^{-f} \tau(\chi)$.
(c) Deduce from (b) that $\tau(\chi)^{e}$ is fixed by the $q^{f}$ th power map, and conclude that $\tau(\chi)^{e} \in \mathbf{F}_{q}\left(\zeta_{e}\right)$.

Table 4. Primes $p=3 \cdot 2^{n}+1$ with $n \leq 750000$ which divide some Fermat number $F_{m}$.

| $n$ | Fermat number $F_{m}$ | Discoverer | Discovered |
| ---: | ---: | ---: | ---: |
| 41 | $F_{38}$ | R. M. Robinson | 1956 |
| 209 | $F_{207}$ | R. M. Robinson | 1956 |
| 157169 | $F_{157167}$ | J. Young | 1995 |
| 213321 | $F_{213319}$ | J. Young | 1996 |
| 303093 | $F_{303088}$ | J. Young | 1998 |
| 382449 | $F_{382447}$ | J. B. Cosgrave \& Y. Gallot | 1999 |

(d) Using (a)-(c), show that
$q$ is an $e$ th power $\bmod p \Longleftrightarrow\left(\tau(\chi)^{e}\right)^{\frac{q^{f}-1}{e}}=1$
$\Longleftrightarrow \tau(\chi)^{e}$ is an $e$ th power in $\mathbf{F}_{q}\left(\zeta_{e}\right)$.
11. (Continuation) Here we consider the cases $e=2$ and $e=3$ which correspond to Gauss's quadratic reciprocity law and Jacobi's cubic reciprocity law.
(a) Let $e=2$, so that the nontrivial character $\chi(\cdot)$ of order 2 can be identified with the Legendre symbol $(\dot{\bar{p}})$. Prove that $\tau_{-1}(\chi)=$ $\chi(-1) \tau_{1}(\chi)$. Using part (a) of the preceding exercise, show that $\tau(\chi)^{2}=\left(\frac{-1}{p}\right) p$, and deduce from part (d) another proof of the law of quadratic reciprocity.
(b) Now suppose $e=3$. One can show that for any $\chi$ as in the preceding exercise, we have $\tau(\chi)^{3}=p \pi$, where $\pi=\frac{L+3 M \sqrt{-3}}{2}$ for certain integers $L, M$ satisfying $L^{2}+27 M^{2}=4 p$ and $L \equiv 1(\bmod 3)(c f$. [Gau86, footnote to Art. 358], [IR90, p. 115]). Assuming this result, deduce another proof of Jacobi's cubic reciprocity law.
12. Give a necessary and sufficient condition in terms of $L$ and $M$ for 6 to be a cubic residue modulo $p$.
13. (Golomb [Gol76])
(a) Suppose $p=3 \cdot 2^{n}+1$ is prime. Show that $p$ divides the $j$ th Fermat number $F_{j}=2^{2^{j}}+1$ for some $j$ if and only if the order of $3(\bmod p)$ is not divisible by 3. Moreover, show that in this case there is exactly one such $j$, and $j<n$.
(b) Prove that if $p=3 \cdot 2^{2 m}+1$ is prime, then the order of 2 modulo $p$ is divisible by 3 , and hence no such primes can divide Fermat numbers. Hint: Show that 2 is not a cubic residue modulo such a prime.
Table 4 lists all primes of the form $3 \cdot 2^{n}+1$ with $n \leq 750000$ which divide a Fermat number.
14. (Kraïtchik, Pellet) Suppose that both $q=2 n+1$ and $p=12 n+7$ are prime. Prove that if $p=L^{\prime 2}+27 M^{\prime 2}$ for integers $L^{\prime}$ and $M^{\prime}$, then $q \mid 2^{p}-1$.

Prove that if both $q=12 n+5$ and $p=72 n+31$ are prime, and $p=L^{\prime 2}+27 M^{\prime 2}$ for integers $L^{\prime}$ and $M^{\prime}$, then $q \mid 2^{p}-1$.

Example: Let $n=18$; then $q=37, p=223=14^{2}+27 \cdot 1^{2}$, and $2^{37}-1=223 \cdot 616318177$.

For other results of this kind see the papers of Fueter [Fue46], Storchi [Sto55] and Golubev [Gol58].
15. Use Kummer's criterion to give another proof that 2 is a cube $\bmod p$ if and only if $2 \mid L$ and $2 \mid M$, and that 3 is a cube $\bmod p$ if and only if $3 \mid M$. Note that these results are less precise than those of Theorems 2.26 and 2.27 . Hint: Before tackling the problem of when 3 is a cube, rewrite the final coefficient of the period polynomial in a form more amenable to computations modulo 3.
16. Prove Theorem 2.37. Use Jacobi's law in the form stated in Theorem 2.28 and the binomial theorem.
17. Prove Theorem 2.38, using Sun's form of Jacobi's reciprocity law.
18. (Lehmer $[$ Leh58]) Let $p \equiv 1(\bmod 3)$ be prime, and suppose $q>3$ is a prime distinct from $p$. Write $4 p=L^{2}+27 M^{2}$. Suppose that $p \equiv \lambda L^{2}(\bmod q)$ for a prime $\lambda$ which can be written in the form $4 \lambda=$ $1+27 m^{2}$ with $q \nmid m$. Show that $q$ is a cube modulo $p$ if and only if $q$ is a cube modulo $\lambda$.

Example (with $\lambda=7, m=1)$ : If $p \equiv 7 L^{2}(\bmod q)$ (equivalently, if $\left.L^{2} \equiv M^{2}(\bmod q)\right)$, then $q$ is a cubic residue modulo $p$ if and only if $q \equiv \pm 1(\bmod 7)$.
19. Let $p \equiv 1(\bmod 3)$ be prime, and write $4 p=L^{2}+27 M^{2}$, where $L \equiv$ $1(\bmod 3)$. For each integer $c$ not divisible by $p$, let $N_{c}$ be the number of ordered pairs $(x, y) \in \mathbf{F}_{p}^{2}$ with $x^{3}+y^{3}=c$.
(a) (Gauss) Show that if $c$ is a cube modulo $p$, then $N_{c}=p-2+L$.
(b) (Chowla, Cowles, \& Cowles [CCC80]) Suppose $c$ is not a cube modulo $p$. Show that $N_{c}=p-2+\frac{1}{2}( \pm 9 M-L)$ and describe how to determine the correct choice of sign.
(c) Deduce that in every case, $\left|N_{c}-(p-2)\right| \leq 2 \sqrt{p}$. This is a special case of a theorem of Hasse known as the Riemann Hypothesis for elliptic curves.
(d) Show that if $p$ is any prime with $p>7$, then every element of $\mathbf{F}_{p}$ is a sum of two cubes. Show, moreover, that if $p>13$, then every element of $\mathbf{F}_{p}$ is a sum of two nonzero cubes.
Hint for ( $b$ ): Give a criterion for $\alpha$ and $c-\alpha$ to be simultaneously cubes in terms of $c^{-1} \alpha$ and $c^{-1} \alpha-1$.

Remark. Leep \& Shapiro [LS89] have shown that if $G$ is a multiplicative subgroup of index 3 in an arbitrary field $F$, then every element of $F$ can be written as a sum of two elements of $G$, unless $\# F=4,7,13$, or 16; see also [BS92].
20. (Gauss [Gau86, footnote to Art. 358], Jacobi [Jac27, Jac69]) Let $p \equiv 1(\bmod 3)$ be prime, say $p=3 f+1$. Write $4 p=L^{2}+27 M^{2}$, where $L \equiv 1(\bmod 3)$. Put

$$
S:=\sum_{\alpha \in \mathbf{F}_{p}^{\times}}\left(\alpha^{3}+1\right)^{2(p-1) / 3} .
$$

(a) Using the binomial theorem, prove that $S=-2-\binom{2 f}{f}$.
(b) Let $g$ be a generator of $\mathbf{F}_{p}^{\times}$and let $\omega$ be the element of $\mathbf{F}_{p}^{\times}$defined by $\omega:=g^{(p-1) / 3}$. Show that, with $a, b$, and $c$ as in Theorem 2.23, we have $S=3 a+3 b \omega^{2}+3 c \omega$.
(c) Check that $\left(\omega^{2}-\omega\right)^{2}=-3$.

In what follows we write " $\sqrt{-3}$ " as an abbreviation for the element $\omega^{2}-\omega \in \mathbf{F}_{p}^{\times}$.
(d) Deduce from (b) and the explicit expressions for $a, b$, and $c$ in Theorem 2.23 that $S=-2+\frac{L+3 M \sqrt{-3}}{2}$.
(e) Conclude that $L+3 M \sqrt{-3}=-2\binom{2 f}{f}$ in $\mathbf{F}_{p}$. Deduce that $L-$ $3 M \sqrt{-3}=0$ in $\mathbf{F}_{p}$.
(f) Show that $L$ is the least absolute remainder of $-\binom{2 f}{f}$ modulo $p$. In other words, $L$ is the unique integer in the interval ( $-p / 2, p / 2$ ) with $L \equiv-\binom{2 f}{f}(\bmod p)$.
Example: Take $p=109=3 \cdot 36+1$. We have $\binom{2 \cdot 36}{36} \equiv 2(\bmod 109)$ and $4 \cdot 109=2^{2}+27 \cdot 4^{2}$.

## Elementary Prime Number Theory, II

Mathematicians have tried in vain to this day to discover some order in the sequence of prime numbers, and we have reason to believe that it is a mystery into which the human mind will never penetrate. - L. Euler

Even before I had begun my more detailed investigations into higher arithmetic, one of my projects was to turn my attention to the decreasing frequency of primes, to which end I counted the primes in several chiliads [intervals of length 1000]... I soon recognized that behind all of its fluctuations, this frequency is on average inversely proportional to the logarithm, so that the number of primes below a given bound $n$ is approximately equal to

$$
\int \frac{d n}{\log n}
$$

where the logarithm is understood to be hyperbolic. - C. F. Gauss

## 1. Introduction

We began our study of prime number theory in Chapter 1 with several different proofs that there are infinitely many primes. In this chapter we turn to the question of how these infinitely many primes are distributed on the real number line. Once again, let $\pi(x)$ denote the number of primes

Table 1. Comparison of $\Delta(x)$ and $1 / \log x$, rounded to the nearest thousandth.

| $x$ | 1000 | 2000 | 3000 | 4000 | 5000 | 6000 | 7000 | 8000 | 9000 | 10000 |
| ---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\Delta(x)$ | .144 | .128 | .122 | .121 | .115 | .117 | .108 | .109 | .118 | 0.107 |
| $\frac{1}{\log x}$ | .145 | .132 | .125 | .121 | .117 | .115 | .113 | .111 | .110 | 0.109 |

$p \leq x$. We would like to understand how quickly and how regularly $\pi(x)$ grows.
1.1. Discovering the prime number theorem. As is the case with much mathematics, the first substantial investigations here were carried out by Gauss. In an 1849 letter to the mathematician and astronomer Encke, Gauss recounted how almost sixty years prior, as a boy of 15 or 16 , he had taken an interest in the function $\pi(x)$.

Gauss's study began with an investigation of what we could term the "local density" of primes near a number $x$. (Some of Gauss's tables have been preserved in [Gau73b, p. 435-443].) Here when we say "local density", what we have in mind is the ratio of the count of primes "near $x$ " with the total number of integers "near $x$ ". Of course this is somewhat vague; Gauss counted primes in intervals of 1000 , which suggests defining

$$
\Delta(x):=\frac{\pi(x+500)-\pi(x-500)}{1000} .
$$

Thus $\Delta(x)$ is the probability of choosing a prime if one samples an integer uniformly at random from the interval $(x-500, x+500]$. Table 1 displays some values of $x$ vs. $\Delta(x)$. From this limited data it appears that $\Delta(x)$ is generally decreasing, albeit somewhat slowly.

But how slowly? To answer this question, Gauss considered the inverse ratio, $\Delta(x)^{-1}$, and discovered empirically that $\Delta(x) \approx 1 / \log x$ (which is also illustrated in Table 1). Since $\Delta(x)$ is the slope of a chord on the graph of $y=\pi(x)$, it is natural to think that one could recover $\pi(x)$ by integrating $1 / \log x$. This suggests that

$$
\begin{equation*}
\pi(x) \approx \int_{2}^{x} \frac{d t}{\log t} \tag{3.1}
\end{equation*}
$$

We use the notation $\operatorname{Li}(x)$ for the integral appearing on the right-hand side of this approximation; it is known as the (Eulerian) logarithmic integral. We refer to (3.1) as the Gauss approximation to $\pi(x)$.

Table 2 compares $\pi(x)$ and $\operatorname{Li}(x)$ for powers of 10 from $10^{3}$ through $10^{13}$. The last column of this table is the most revealing. It suggests that for larger and larger values of $x$, the Gauss approximation very quickly approaches $100 \%$ accuracy. In other words, it seems that the following is

Table 2. Comparison of $\pi(x)$ and $\operatorname{Li}(x)$, where $\operatorname{Li}(x)$ is rounded to the nearest integer. The last column gives the percentage error, computed as $|\operatorname{Li}(x)-\pi(x)| / \pi(x)$.

| $x$ | $\pi(x)$ | $\operatorname{Li}(x)$ | $\operatorname{Li}(x)-\pi(x)$ | $\%$ error |
| ---: | ---: | ---: | ---: | ---: |
| $10^{3}$ | 168 | 177 | 9 | $5.4 \%$ |
| $10^{4}$ | 1229 | 1245 | 16 | $1.3 \%$ |
| $10^{5}$ | 9,592 | 9,629 | 37 | $3.8 \times 10^{-1} \%$ |
| $10^{6}$ | 78,498 | 78,627 | 129 | $1.6 \times 10^{-1} \%$ |
| $10^{7}$ | 664,579 | 664,917 | 338 | $5.1 \times 10^{-2} \%$ |
| $10^{8}$ | $5,761,455$ | $5,762,208$ | 753 | $1.3 \times 10^{-2} \%$ |
| $10^{9}$ | $50,847,534$ | $50,849,234$ | 1,700 | $3.3 \times 10^{-3} \%$ |
| $10^{10}$ | $455,052,512$ | $455,055,614$ | 3,102 | $6.8 \times 10^{-4} \%$ |
| $10^{11}$ | $4,118,054,813$ | $4,118,066,400$ | 11,587 | $2.8 \times 10^{-4} \%$ |
| $10^{12}$ | $37,607,912,018$ | $37,607,950,280$ | 38,262 | $1.0 \times 10^{-4} \%$ |
| $10^{13}$ | $346,065,536,839$ | $346,065,645,809$ | 108,970 | $3.2 \times 10^{-5} \%$ |

true:
$\star$ Theorem 3.1 (Prime number theorem). $\pi(x) \sim \operatorname{Li}(x)$ as $x \rightarrow \infty$.
In 1859, Riemann outlined a strategy for proving Theorem 3.1 based on viewing the function $\zeta(s)=\sum_{n=1}^{\infty} n^{-s}$, introduced by Euler, as a function of a complex variable $s$. But it took until 1896 for complex analysis to mature to the point where Riemann's outline could be filled in; this was done independently by Hadamard and de la Vallée-Poussin. There are still no simple proofs of Theorem 3.1, although there are short proofs which require only a modicum of familiarity with complex analysis (see, e.g., [Zag97]). In Chapter 7, we will give a (long) proof of the prime number theorem completely independent of the theory of complex variables.
1.2. An alternative formulation of the prime number theorem. The prime number theorem is often stated in the following simpler form:
$\star$ Theorem 3.2 (Prime number theorem, alternative form). As $x \rightarrow \infty$, $\pi(x) \sim x / \log x$.

It is not difficult to show that Theorems 3.1 and 3.2 are equivalent: If we integrate $1 / \log t$ by parts, we find that

$$
\begin{align*}
\operatorname{Li}(x)=\int_{2}^{x} \frac{d t}{\log t} & =\left.\frac{t}{\log t}\right|_{2} ^{x}+\int_{2}^{x} \frac{d t}{(\log t)^{2}} \\
& =\frac{x}{\log x}-\frac{2}{\log 2}+\int_{2}^{x} \frac{d t}{(\log t)^{2}} . \tag{3.2}
\end{align*}
$$

Moreover, the final integral is $o(\operatorname{Li}(x))$. Indeed, by L'Hôpital's rule,

$$
\lim _{x \rightarrow \infty} \frac{\int_{2}^{x} \frac{d t}{(\log t)^{2}}}{\operatorname{Li}(x)}=\lim _{x \rightarrow \infty} \frac{1 /(\log x)^{2}}{1 / \log x}=0
$$

Hence $\operatorname{Li}(x) \sim x / \log x$, from which the equivalence of Theorems 3.1 and 3.2 follows.
1.3. What happens now? Since the prime number theorem is not proved until Chapter 7, what is left for us to do here? Prior to the proof of Theorem 3.1, several estimates for quantities related to $\pi(x)$ were obtained by Chebyshev, Mertens, and others. For many applications, these are more than sufficient; the prime number theorem itself is not required. In particular, this comment applies to our treatment of sieve methods in Chapter 6. Moreover, these estimates are necessary preliminaries for our eventual proof of the prime number theorem. We devote most of this chapter to a discussion of these results and their charming, elementary proofs.

In the final section we revisit Gauss's heuristic for the prime number theorem. We explain how Gauss's observation that the "local density" of the primes near $x$ is $\approx 1 / \log x$ suggests many other statements about primes. For example, we show how Gauss's idea can be used to formulate a plausible prediction of the number of twin prime pairs up to $x$.

## 2. The set of prime numbers has density zero

After a moment's reflection on the definitions, most intelligent laymen can convince themselves that the prime numbers account for at most half of the natural numbers. Indeed, one of the first facts people tend to notice about the primes is that every prime number $p>2$ is odd. A small elaboration on this trivial observation permits one to establish the following:

Theorem 3.3. $\pi(x) / x \rightarrow 0$ as $x \rightarrow \infty$. That is, the set of primes has asymptotic density zero.

Proof. Let $q$ be any (fixed) natural number. Then every prime $p$ that does not divide $q$ belongs to one of the $\varphi(q)$ invertible residue classes modulo $q$. The number of natural numbers $n \leq x$ which fall in a given residue class modulo $q$ is at most $1+x / q$, and so the number of $n \leq x$ which are coprime to $q$ is at most $\varphi(q)+x \varphi(q) / q$. Since only finitely many primes $p$ divide $q$, this shows that

$$
\pi(x) \leq(\varphi(q) / q+o(1)) x \quad(x \rightarrow \infty)
$$

Theorem 3.3 will follow if we can show that $\varphi(q) / q$ can be made arbitrarily small. For each $z>0$, put $q:=q_{z}=\prod_{p \leq z} p$. From (1.4), we
have

$$
\frac{\varphi\left(q_{z}\right)}{q_{z}}=\prod_{p \leq z}\left(1-\frac{1}{p}\right) \leq \exp \left(-\sum_{p \leq z} \frac{1}{p}\right) .
$$

Since $\sum_{p} p^{-1}$ diverges, it follows that $\varphi\left(q_{z}\right) / q_{z} \rightarrow 0$ as $z \rightarrow \infty$.
It is remarkable that a result asserting that there are not too few primes (namely, that $\sum_{p} p^{-1}$ diverges) is used here to show that there are not too many primes (Theorem 3.3). Actually, if we assume (contrary to fact) that $\sum_{p} p^{-1}$ converges, it is also easy to show that $\pi(x) / x \rightarrow 0$; see Exercise 1.

## 3. Three theorems of Chebyshev

[Chebyshev] was the only man ever able to cope with the refractory character and erratic flow of prime numbers and to confine the stream of their progression with algebraic limits, building up, if I may so say, banks on either side which that stream, devious and irregular as are its windings, can never overflow. - J. J. Sylvester

In 1851 and 1852, Chebyshev published two important papers [Che51, Che52] on the behavior of $\pi(x)$. We shall focus our attention on three of his results:

Theorem 3.4. If $\frac{\pi(x)}{x / \log x}$ tends to a limit as $x \rightarrow \infty$, then that limit is 1 .
Theorem 3.5. There exist positive constants $c_{1}, c_{2}$ and a real number $x_{0}$ so that

$$
c_{1} \frac{x}{\log x} \leq \pi(x) \leq c_{2} \frac{x}{\log x} \quad\left(\text { whenever } x>x_{0}\right) .
$$

Theorem 3.5 shows that the prime number theorem at least predicts the correct order of magnitude of $\pi(x)$. Theorem 3.4 shows that if $\pi(x)$ behaves regularly enough that $\pi(x) \sim c x / \log x$ for some constant $c$, then the prime number theorem holds. (For a more general result of the same character as Theorem 3.4, see Exercises 28 and 29.)

Theorem 3.6 (Bertrand's postulate). For all sufficiently large $x$, there is a prime in the interval $(x, 2 x]$.

Actually Bertrand conjectured, and Chebyshev proved, that the conclusion of Theorem 3.6 is valid for every real $x \geq 1$. This follows from the argument presented below after a finite computation; cf. Exercises 12-13.

Before proving these results, it is convenient to introduce certain auxiliary functions. Put

$$
\begin{equation*}
\theta(x):=\sum_{p \leq x} \log p, \quad \psi(x):=\sum_{n=1}^{\infty} \theta\left(x^{1 / n}\right) . \tag{3.3}
\end{equation*}
$$

The sum defining $\psi$ appears to be infinite, but is morally finite since $\theta\left(x^{1 / n}\right)$ vanishes once $x^{1 / n}<2$. The functions $\psi$ and $\theta$ turn out to be better-behaved and easier to study than $\pi(x)$. Fortunately, estimates for $\pi(x)$ can be easily deduced from estimates for either $\theta$ or $\psi$ : By partial summation,

$$
\theta(x)=\pi(x) \log x-\int_{2}^{x} \frac{\pi(t)}{t} d t
$$

Because $\pi(t) / t=o(1)$ (Theorem 3.3), we have $\int_{2}^{x} \pi(t) / t d t=o(x)$, whence

$$
\theta(x)=\pi(x) \log x+o(x)
$$

and

$$
\begin{equation*}
\frac{\theta(x)}{x}=\frac{\pi(x)}{x / \log x}+o(1) . \tag{3.4}
\end{equation*}
$$

The analogue of (3.4) holds with $\psi$ in place of $\theta$, because the difference between $\psi$ and $\theta$ is quite small: Indeed, write

$$
\begin{equation*}
\psi(x)-\theta(x)=\theta\left(x^{1 / 2}\right)+\theta\left(x^{1 / 3}\right)+\cdots . \tag{3.5}
\end{equation*}
$$

As observed above, $\theta\left(x^{1 / n}\right)$ vanishes whenever $x^{1 / n}<2$, i.e., once $n>$ $\log x / \log 2$. Consequently, only $O(\log x)$ of the terms on the right of (3.5) are nonzero. Because $\theta(t) \leq \sum_{n \leq t} \log t \leq t \log t$ trivially,

$$
\begin{equation*}
\psi(x)-\theta(x) \ll x^{1 / 2} \log x+\left(x^{1 / 3} \log x\right) \log x \ll x^{1 / 2} \log x \tag{3.6}
\end{equation*}
$$

Thus replacing $\theta$ with $\psi$ in equation (3.4) results in an extra error term of $O\left((\log x) x^{-1 / 2}\right)$, which can be absorbed into the existing $o(1)$ error term. Thus we have proved:

Proposition 3.7. As $x \rightarrow \infty$, we have both

$$
\begin{align*}
\frac{\theta(x)}{x} & =\frac{\pi(x)}{x / \log x}+o(1)  \tag{3.7}\\
\frac{\psi(x)}{x} & =\frac{\pi(x)}{x / \log x}+o(1) \tag{3.8}
\end{align*}
$$

This has the following useful consequence:
Corollary 3.8. If any of $\frac{\theta(x)}{x}$, $\frac{\psi(x)}{x}$, or $\frac{\pi(x)}{x / \log x}$ tends to a limit as $x \rightarrow \infty$, then all of them do, and the limit in each case is the same. In particular, the prime number theorem is equivalent to the estimate $\theta(x) \sim x$ and to the estimate $\psi(x) \sim x$.

Indeed, (3.7) and (3.8) together imply that

$$
\liminf _{x \rightarrow \infty} \frac{\theta(x)}{x}=\liminf _{x \rightarrow \infty} \frac{\pi(x)}{x / \log x}=\liminf _{x \rightarrow \infty} \frac{\psi(x)}{x},
$$

and similarly for the lim sup.
The definition of $\psi$ given above is useful for making comparisons with $\theta$, but it masks the arithmetic information that $\psi$ encodes. To get at this, observe that for any fixed positive integer $k$,

$$
\theta\left(x^{1 / k}\right)=\sum_{p \leq x^{1 / k}} \log p=\sum_{p^{k} \leq x} \log p .
$$

Hence

$$
\begin{equation*}
\psi(x)=\theta(x)+\theta\left(x^{1 / 2}\right)+\cdots=\sum_{p^{k} \leq x} \log p, \tag{3.9}
\end{equation*}
$$

where the final sum is over all pairs $(p, k)$ where $p$ is prime, $k$ is a positive integer and $p^{k} \leq x$. Define the von Mangoldt function $\Lambda(n)$ by

$$
\Lambda(n):= \begin{cases}\log p & \text { if } n=p^{k} \text { is a prime power }, \\ 0 & \text { otherwise }\end{cases}
$$

The fundamental theorem of arithmetic assures us that $\Lambda$ is well-defined, and from equation (3.9) we can read off the identity

$$
\psi(x)=\sum_{n \leq x} \Lambda(n) .
$$

Lemma 3.9. For every positive integer $n$,

$$
\sum_{d \mid n} \Lambda(n)=\log n
$$

Proof. Write $n=\prod_{p \mid n} p^{e_{p}}$. Then

$$
\begin{aligned}
\sum_{d \mid n} \Lambda(d)=\sum_{p^{k} \mid n} \log p & =\sum_{p \mid n} \sum_{k=1}^{e_{p}} \log p \\
& =\sum_{p \mid n} e_{p} \log p=\sum_{p \mid n} \log p^{e_{p}}=\log \left(\prod_{p \mid n} p^{e_{p}}\right)=\log n .
\end{aligned}
$$

Following Chebyshev, we now set $T(x):=\sum_{n \leq x} \log n$.
Lemma 3.10. For $x \geq 2$, we have

$$
\begin{equation*}
T(x)=x \log x-x+O(\log x) \tag{3.10}
\end{equation*}
$$

Proof. Since $\log t$ is increasing for $t>0$, we have $\log n \leq \int_{n}^{n+1} \log t d t \leq$ $\log (n+1)$ for each natural number $n$. So

$$
\sum_{n \leq x} \log n \leq \int_{1}^{\lfloor x\rfloor+1} \log t d t=(\lfloor x\rfloor+1) \log (\lfloor x\rfloor+1)-(\lfloor x\rfloor+1)+1
$$

and

$$
\sum_{n \leq x} \log n=\sum_{2 \leq n \leq x} \log n \geq \int_{1}^{\lfloor x\rfloor} \log t d t=\lfloor x\rfloor \log \lfloor x\rfloor-\lfloor x\rfloor+1 .
$$

Both the upper and lower bounds are $x \log x-x+O(\log x)$, and so the lemma follows.

The link between $T(x)$ and prime number theory is given by the following result, which is the fundamental tool in the proofs of Theorems 3.4-3.6.

Lemma 3.11. For every $x>0$, we have $T(x)=\sum_{n \leq x} \psi(x / n)$.
Proof. Observe that

$$
\begin{aligned}
\sum_{n \leq x} \psi(x / n)=\sum_{n \leq x} \sum_{m \leq x / n} \Lambda(m) & =\sum_{n m \leq x} \Lambda(m) \\
& =\sum_{N \leq x} \sum_{m \mid N} \Lambda(m)=\sum_{N \leq x} \log N=T(x)
\end{aligned}
$$

3.1. Proof of Theorem 3.4. We begin with a plausibility argument for the prime number theorem: From Lemma 3.11 and (3.10),

$$
\begin{equation*}
\sum_{n \leq x} \psi(x / n) \sim x \log x \quad(x \rightarrow \infty) \tag{3.11}
\end{equation*}
$$

This is the same estimate one would obtain if the terms on the left of (3.11) were " $x / n$ " instead of " $\psi(x / n)$ ", which can be considered evidence for the prime number theorem in the form $\psi(x) \sim x$.

This idea can be used to prove the following proposition, which in view of Proposition 3.7 implies Theorem 3.4.

Proposition 3.12. We have

$$
\liminf _{x \rightarrow \infty} \frac{\psi(x)}{x} \leq 1 \leq \limsup _{x \rightarrow \infty} \frac{\psi(x)}{x}
$$

Proof. Put $c:=\liminf _{x \rightarrow \infty} \psi(x) / x$ and $C:=\limsup _{x \rightarrow \infty} \psi(x) / x$. Then $\psi(x) \geq c x+g(x)$ for a function $g(x)$ satisfying $g(x)=o(x)$. Hence

$$
\begin{align*}
\sum_{n \leq x} \psi(x / n) & \geq c x \sum_{n \leq x} n^{-1}+\sum_{n \leq x} g(x / n) \\
& =c x \log x+o(x \log x)+\sum_{n \leq x} g(x / n) \tag{3.12}
\end{align*}
$$

We claim that the final summand can be absorbed into the error term $o(x \log x)$. This implies that $\sum_{n \leq x} \psi(x / n) \geq(c+o(1)) x \log x$, which (by (3.11)) implies $c \leq 1$. A similar argument, with $c$ replaced by $C$, shows that $C \geq 1$.

To prove the claim about $\sum g(x / n)$, let $\epsilon>0$ be given and choose $N$ so large that $|g(t)| t^{-1}<\epsilon / 2$ whenever $t>N$. Let $M$ be an upper bound for $|g|$ on $[1, N]$. Then

$$
\begin{aligned}
\left|\sum_{n \leq x} g(x / n)\right| & \leq \sum_{\substack{n \leq x \\
x / n \leq N}}|g(x / n)|+\sum_{\substack{n \leq x \\
x / n>N}}|g(x / n)| \\
& \leq M x+\frac{\epsilon}{2} x \sum_{n \leq x} n^{-1}<\epsilon x \log x
\end{aligned}
$$

for sufficiently large $x$.
3.2. Proof of Theorem 3.5. Suppose $x \geq 4$. By Lemma 3.10,

$$
\begin{aligned}
T(x)-2 T(x / 2) & =x \log x-x+O(\log x)-2\left(\frac{x}{2} \log \frac{x}{2}-\frac{x}{2}+O\left(\log \frac{x}{2}\right)\right) \\
& =x \log 2+O(\log x)
\end{aligned}
$$

On the other hand, Lemma 3.11 implies that

$$
\begin{aligned}
T(x)-2 T(x / 2) & =\sum_{n \leq x} \psi(x / n)-\sum_{n \leq x} 2 \psi(x / 2 n) \\
& =\sum_{n \geq 1}(-1)^{n-1} \psi(x / n)=\psi(x)-\psi(x / 2)+\cdots .
\end{aligned}
$$

Since $\psi$ is an increasing function, this is an alternating series of decreasing terms. It follows that for any even $k$,

$$
\begin{equation*}
T(x)-2 T(x / 2) \geq \psi(x)-\psi(x / 2)+\cdots+\psi(x /(k-1))-\psi(x / k), \tag{3.13}
\end{equation*}
$$

while for any odd $k$,

$$
\begin{equation*}
T(x)-2 T(x / 2) \leq \psi(x)-\psi(x / 2)+\cdots-\psi(x /(k-1))+\psi(x / k) . \tag{3.14}
\end{equation*}
$$

Taking $k=1$ in (3.14) gives the lower bound

$$
\begin{equation*}
\psi(x) \geq T(x)-2 T(x / 2)=x \log 2+O(\log x) . \tag{3.15}
\end{equation*}
$$

Getting an upper bound on $\psi(x)$ is a tad bit trickier. First take $k=2$ in (3.13) to find that

$$
\psi(x)-\psi(x / 2) \leq T(x)-2 T(x / 2)=x \log 2+O(\log x)
$$

Now let $k$ be the positive integer for which $x / 2^{k-1} \geq 4>x / 2^{k}$. For each $1 \leq j \leq k$,

$$
\psi\left(x / 2^{j-1}\right)-\psi\left(x / 2^{j}\right) \leq \frac{x}{2^{j-1}} \log 2+O\left(\log \frac{x}{2^{j-1}}\right)=\frac{x}{2^{j-1}} \log 2+O(\log x)
$$

Summing these inequalities for $1 \leq j \leq k$, we have (noting that $k \ll \log x$ )

$$
\begin{aligned}
\psi(x)-\psi\left(x / 2^{k}\right) & \leq x \log 2\left(1+\frac{1}{2}+\cdots+\frac{1}{2^{k-1}}\right)+O((\log x)(\log x)) \\
& \leq 2 x \log 2+O\left((\log x)^{2}\right)
\end{aligned}
$$

Thus

$$
\begin{equation*}
\psi(x) \leq 2 x \log 2+O\left((\log x)^{2}\right)+\psi(4) \leq 2 x \log 2+O\left((\log x)^{2}\right) . \tag{3.16}
\end{equation*}
$$

Collecting our upper and lower bounds on $\psi(x)$, we have proved:
Proposition 3.13. For $x \geq 4$, we have

$$
\begin{equation*}
x \log 2+O(\log x) \leq \psi(x) \leq 2 x \log 2+O\left((\log x)^{2}\right) . \tag{3.17}
\end{equation*}
$$

From Propositions 3.7 and 3.13, we obtain Theorem 3.5 for any constants $c_{1}$ and $c_{2}$ satisfying $c_{1}<\log 2$ and $c_{2}>2 \log 2$. Since $\frac{2 \log 2}{\log 2}=2$, this has the following corollary: For each fixed $\epsilon>0$, there is a prime in the interval $[x,(2+\epsilon) x]$ for all $x>x_{0}(\epsilon)$. Said differently, we are an $\epsilon$ away from a proof of Bertrand's postulate!
3.3. Proof of Bertrand's postulate. Obviously, if we can produce a nonvanishing sum over the primes $p \in(x, 2 x]$, then there must be a prime in $(x, 2 x]$. In particular, Bertrand's postulate will follow if we show that

$$
\theta(2 x)-\theta(x)=\sum_{x<p \leq 2 x} \log p>0
$$

for large enough $x$. We will establish this by first estimating $\psi(2 x)-\psi(x)$ from below, and then using (3.6) to translate that estimate into a lower bound on $\theta(2 x)-\theta(x)$.

Here one's first instinct is perhaps to take $k=2$ in (3.13), as this immediately gives us a bound on $\psi(x)-\psi(x / 2)$, namely

$$
\psi(x)-\psi(x / 2) \leq T(x)-2 T(x / 2) .
$$

Unfortunately, the inequality is going the wrong way for our purposes. So instead we take $k=3$ in (3.14); this gives us that

$$
\begin{equation*}
\psi(x)-\psi(x / 2)+\psi(x / 3) \geq T(x)-2 T(x / 2)=x \log 2+O(\log x) . \tag{3.18}
\end{equation*}
$$

This inequality is going the right way but has the extra term $\psi(x / 3)$. However, from (3.16),

$$
\begin{equation*}
\psi(x / 3) \leq \frac{2 \log 2}{3} x+O\left((\log x)^{2}\right) \tag{3.19}
\end{equation*}
$$

which in conjunction with (3.18) implies that

$$
\psi(x)-\psi(x / 2) \geq x \frac{\log 2}{3}+O\left((\log x)^{2}\right)
$$

Invoking (3.6), we obtain the lower bound

$$
\begin{equation*}
\theta(x)-\theta(x / 2) \geq x \frac{\log 2}{3}+O\left(x^{1 / 2} \log x\right) \quad(x \rightarrow \infty) . \tag{3.20}
\end{equation*}
$$

Theorem 3.6 is now immediate, since the right-hand side of (3.20) is positive for large $x$.

In fact, (3.20) yields a lower bound for $\pi(x)-\pi(x / 2)$ of the same order of magnitude as the lower bound for $\pi(x)$ in Theorem 3.5. Indeed,

$$
\theta(x)-\theta(x / 2)=\sum_{x / 2<p \leq x} \log p \leq \log x(\pi(x)-\pi(x / 2)),
$$

so that by (3.20),

$$
\begin{equation*}
\pi(x)-\pi(x / 2) \geq \frac{\log 2}{3} \frac{x}{\log x}+O\left(x^{1 / 2}\right)=\left(\frac{\log 2}{3}+o(1)\right) \frac{x}{\log x} \quad(x \rightarrow \infty) \tag{3.21}
\end{equation*}
$$

This proof of Bertrand's postulate is due to Ramanujan [Ram19].

## 4. The work of Mertens

By 1737, Euler was aware not only of the divergence of $\sum_{p} p^{-1}$, but had assigned the infinite sum the value $\log \log \infty[E u l 37$, Theorema 19], showing that he possessed an inkling as to the rate of growth of the partial sums. In Gauss's Nachlass [Gau73c, pp. 11-16] one can find the more precise assertion that

$$
" 1+\frac{1}{2}+\frac{1}{3}+\frac{1}{5}+\frac{1}{7}+\frac{1}{11}+\cdots+\frac{1}{x}=(\text { for } x \text { infinite }) l l x+V . "
$$

Gauss writes that he suspects $V$ to be a constant near 1.266. It seems reasonable to read this as the conjecture that

$$
\sum_{p \leq x} \frac{1}{p}=\log \log x+V-1+o(1)
$$

Gauss also claims that

$$
" \frac{2}{1} \cdot \frac{3}{2} \cdot \frac{5}{4} \cdots \frac{x}{x-1}=(x \text { inf }) a . l x "
$$

for a constant $a \approx 1.874$, which we can read as the conjecture that

$$
\prod_{p \leq x}(1-1 / p)^{-1} \sim a \log x .
$$

Mertens observed [Mer74] that Chebyshev's results could be used to obtain precise estimates for both $\sum_{p \leq x} 1 / p$ and $\prod_{p \leq x}(1-1 / p)$. His results vindicate Gauss's claims, apart from small inaccuracies in the numerical values of the constants; the correct values are $V=1.2614972 \ldots$ and $a=$ 1.7810724....
4.1. Mertens' first theorem. We begin by considering the weighted sum $A(x):=\sum_{p \leq x} \log p / p$. From estimates for $A(x)$, results on $\sum_{p \leq x} 1 / p$ follow by partial summation, and these in turn easily yield theorems about $\prod_{p \leq x}(1-1 / p)$.

Observe that the function $T(x)$ introduced in $\S 3$ can be written in the form

$$
T(x)=\sum_{n \leq x} \log n=\sum_{n \leq x} \sum_{d \mid n} \Lambda(d)=\sum_{\substack{d \leq x}} \sum_{\substack{n \leq x \\ d \mid n}} \Lambda(d)=\sum_{d \leq x} \Lambda(d)\left\lfloor\frac{x}{d}\right\rfloor .
$$

If we drop the greatest integer sign, then the error incurred in the sum is $\ll \sum_{d \leq x} \Lambda(d)=\psi(x) \ll x$ by (3.17). Now substituting in the estimate $T(x)=x \log x+O(x)$ furnished by Lemma 3.10 and dividing by $x$, we are led to the important result that

$$
\begin{equation*}
\sum_{d \leq x} \frac{\Lambda(d)}{d}=\log x+O(1) \tag{3.22}
\end{equation*}
$$

Observe that

$$
\begin{equation*}
\sum_{d \leq x} \frac{\Lambda(d)}{d}=\sum_{p^{k} \leq x} \frac{\log p}{p^{k}} \tag{3.23}
\end{equation*}
$$

So if it were not for the terms corresponding to prime powers $p^{k}$ with $k \geq$ 2 , then (3.22) would be an estimate for $A(x)$. But these nuisance terms contribute a bounded amount:

$$
\begin{align*}
\sum_{\substack{p^{k} \leq x \\
k \geq 2}} \frac{\log p}{p^{k}} & \leq \sum_{p \leq x} \log p \sum_{k=2}^{\infty} p^{-k} \\
& =\sum_{p \leq x} \frac{\log p}{p(p-1)} \leq \sum_{2 \leq n \leq x} \frac{\log n}{n(n-1)}=O(1) . \tag{3.24}
\end{align*}
$$

Combining (3.22), (3.23) and (3.24), we obtain that (for $x \geq 1$ )

$$
\begin{equation*}
A(x)=\sum_{p \leq x} \frac{\log p}{p}=\log x+O(1) \tag{3.25}
\end{equation*}
$$

Theorem 3.14 (Mertens' first theorem). As $x \rightarrow \infty$, we have $\sum_{p \leq x} p^{-1}=$ $\log \log x+B_{1}+O(1 / \log x)$ for a constant $B_{1}$. Here $B_{1}=1-\log \log 2+$ $\int_{2}^{\infty}(A(t)-\log t) /\left(t(\log t)^{2}\right) d t$.

Proof. By partial summation,

$$
\sum_{p \leq x} \frac{1}{p}=\sum_{p \leq x} \frac{\log p}{p} \frac{1}{\log p}=\frac{A(x)}{\log x}+\int_{2}^{x} \frac{A(t)}{t(\log t)^{2}} d t
$$

From (3.25) we have that $A(x) / \log x=1+O(1 / \log x)$. To estimate the integral, we write $A(t)=\log t+(A(t)-\log t)$, so that

$$
\begin{aligned}
\int_{2}^{x} \frac{A(t)}{t(\log t)^{2}} d t & =\int_{2}^{x} \frac{1}{t \log t} d t+\int_{2}^{x} \frac{A(t)-\log t}{t(\log t)^{2}} d t \\
& =\log \log x-\log \log 2+\int_{2}^{x} \frac{A(t)-\log t}{t(\log t)^{2}} d t
\end{aligned}
$$

Since $A(t)-\log t$ is bounded, the integral $I:=\int_{2}^{\infty} \frac{A(t)-\log t}{t(\log t)^{2}} d t$ converges absolutely. Moreover,

$$
I-\int_{2}^{x} \frac{A(t)-\log t}{t(\log t)^{2}} d t \ll \int_{x}^{\infty} \frac{d t}{t(\log t)^{2}}=\frac{1}{\log x}
$$

Piecing everything together yields the theorem.
4.2. Mertens' second theorem. The second theorem of Mertens, which is usually the result intended when one sees references to Mertens' theorem in the literature, governs the behavior of the product $\prod_{p \leq x}(1-1 / p)$.
Theorem 3.15. There is an absolute constant $C$ for which $\prod_{p \leq x}(1-1 / p)=$ $e^{-C} / \log x+O\left(1 /(\log x)^{2}\right)$ as $x \rightarrow \infty$. Explicitly, $C=B_{1}+B_{2}$, where $B_{1}$ is the constant of Theorem 3.14 and $B_{2}:=\sum_{p} \sum_{k=2}^{\infty}\left(k p^{k}\right)^{-1}$.

Proof. Let $P_{x}:=\prod_{p \leq x}(1-1 / p)$. Since $\log (1-1 / p)=-\sum_{k \geq 1}\left(k p^{k}\right)^{-1}$,

$$
\log P_{x}=-\sum_{p \leq x} \frac{1}{p}-\sum_{p \leq x} \sum_{k=2}^{\infty} \frac{1}{k p^{k}} .
$$

Since

$$
\sum_{k=2}^{\infty} \frac{1}{k p^{k}} \leq \frac{1}{2} \sum_{k=2}^{\infty} \frac{1}{p^{k}}=\frac{1}{2 p(p-1)} \leq \frac{1}{p^{2}}
$$

the infinite sum $\sum_{p} \sum_{k=2}^{\infty}\left(k p^{k}\right)^{-1}$ converges absolutely, to $B_{2}$, say. Moreover, $B_{2}-\sum_{p \leq x} \sum_{k=2}^{\infty}\left(k p^{k}\right)^{-1} \leq \sum_{p>x} p^{-2} \ll x^{-1}$. Hence

$$
\begin{aligned}
\log P_{x} & =-\log \log x-B_{1}+O(1 / \log x)-B_{2}+O(1 / x) \\
& =-\log \log x-B_{1}-B_{2}+O(1 / \log x) .
\end{aligned}
$$



Figure 1.

Exponentiating, we find that

$$
\prod_{p \leq x}\left(1-\frac{1}{p}\right)=\frac{\exp \left(-\left(B_{1}+B_{2}\right)\right)}{\log x} \exp (O(1 / \log x))
$$

and the result follows with $C=B_{1}+B_{2}$.

In the remainder of this section we show that the constant $C$ of Theorem 3.15 admits a much more pleasant description.

Lemma 3.16 (Euler). For $x \geq 1$, we have $\sum_{n \leq x} n^{-1}=\log x+\gamma+O(1 / x)$, where $\gamma$ is an absolute constant.

The constant $\gamma=0.57721566490153286061 \ldots$ is known as the EulerMascheroni constant.

Proof. Let $r_{n}=n^{-1}-\int_{n}^{n+1} t^{-1} d t$. Then $r_{n}$ is the area of that part of the rectangle $[n, n+1] \times[0,1 / n]$ that lies above the graph of $y=1 / x$. From Figure 1 it is clear that $\sum_{n \geq 1} r_{n}$ converges to a number $\gamma$ less than 1 . For each natural number $N$, we have

$$
\sum_{n \leq N} r_{n}=\sum_{n \leq N} \frac{1}{n}-\int_{1}^{N+1} \frac{d t}{t}
$$

Thus

$$
\begin{aligned}
\sum_{n \leq N} \frac{1}{n} & =\int_{1}^{N+1} \frac{d t}{t}+\sum_{n \leq N} r_{n} \\
& =\log (N+1)+\gamma-\sum_{n=N+1}^{\infty} r_{n}
\end{aligned}
$$

From Figure 1 it is clear that $\sum_{n=N+1}^{\infty} r_{n} \leq(N+1)^{-1}$. So, taking $N=\lfloor x\rfloor$, we deduce that for $x \geq 1$, we have $\sum_{n \leq x} n^{-1}=\log (\lfloor x\rfloor+1)+\gamma+O(1 / x)$. Since $\log (\lfloor x\rfloor+1)=\log x+O(1 / x)$ for $x \geq 1$, the lemma follows.

Theorem 3.17. In the notation of Theorems 3.14 and 3.15, we have $C=$ $B_{1}+B_{2}=\gamma$, where $\gamma$ is the Euler-Mascheroni constant.

Proof. For real $s>1$, we let $\zeta(s)=\sum_{n=1}^{\infty} n^{-s}$ be the Euler-Riemann zeta function (introduced in Chapter 1, §4) and we let $Z(s):=\sum_{p} p^{-s}$. Put $F(s):=\log \zeta(s)-Z(s)$. Since $\zeta(s)=\prod_{p}\left(1-p^{-s}\right)^{-1}$, a short calculation shows that $F(s)=\sum_{k, p}^{\prime}\left(k p^{k s}\right)^{-1}$, where the ${ }^{\prime}$ indicates that the sum is over primes $p$ and integers $k \geq 2$. This series for $F(s)$ converges uniformly on each compact subset of $(1 / 2, \infty)$, and so $F(s) \rightarrow \sum^{\prime}\left(k p^{k}\right)^{-1}=B_{2}$ as $s \downarrow 1$.

We now derive an alternative representation for $F(s)$ which will make visible that $F(s) \rightarrow \gamma-B_{1}$ as $s \downarrow 1$, thus proving that $\gamma=B_{1}+B_{2}$. We start by noting that since $t^{-s}$ is decreasing for $t>0$ (for each fixed $s>1$ ),

$$
\frac{1}{s-1}=\int_{1}^{\infty} t^{-s} d t \leq \zeta(s) \leq 1+\int_{1}^{\infty} t^{-s} d t=1+\frac{1}{s-1}
$$

so that $0 \leq \zeta(s)-(s-1)^{-1} \leq 1$. Hence $\log \zeta(s)=\log \left((s-1)^{-1}\right)+O(s-1)$. Since $1-e^{-(s-1)}=(s-1)(1+O(s-1))$, it follows that $\log \left(1-e^{-(s-1)}\right)=$ $\log (s-1)+O(s-1)$, and so

$$
\begin{align*}
\log \zeta(s) & =-\log \left(1-e^{-(s-1)}\right)+O(s-1) \\
& =\sum_{n=1}^{\infty} e^{-(s-1) n} n^{-1}+O(s-1) . \tag{3.26}
\end{align*}
$$

With $H(x):=\sum_{n \leq x} n^{-1}$, the sum in (3.26) is $\int_{0}^{\infty} e^{-(s-1) t} d H(t)$, which (after a short calculation) shows that

$$
\log \zeta(s)=(s-1) \int_{0}^{\infty} H(t) e^{-(s-1) t} d t+O(s-1)
$$

Let $P(x):=\sum_{p \leq x} p^{-1}$. Another application of partial summation shows that

$$
Z(s)=(s-1) \int_{1}^{\infty} t^{-s} P(t) d t=(s-1) \int_{0}^{\infty} e^{-(s-1) t} P\left(e^{t}\right) d t
$$

Theorem 3.14 implies that $P\left(e^{t}\right)=\log t+B_{1}+O\left((t+1)^{-1}\right)$ for $t \geq 0$ and Theorem 3.16 gives us that $H(t)=\log t+\gamma+O\left((t+1)^{-1}\right)$ for $t \geq 1$. So

$$
\begin{aligned}
F(s) & =\log \zeta(s)-Z(s) \\
& =(s-1) \int_{0}^{\infty} e^{-(s-1) t}\left(H(t)-P\left(e^{t}\right)\right) d t+O(s-1) \\
& =(s-1) \int_{0}^{\infty} e^{-(s-1) t}\left(\gamma-B_{1}+O\left(\frac{1}{t+1}\right)\right) d t+O(s-1) .
\end{aligned}
$$

Here the main term is

$$
(s-1) \int_{0}^{\infty} e^{-(s-1) t}\left(\gamma-B_{1}\right) d t=\gamma-B_{1}
$$

and the error term is

$$
\ll(s-1)+(s-1) \int_{0}^{\infty} \frac{e^{-(s-1) t}}{t+1} d t .
$$

Splitting this last integral at $t=(s-1)^{-1}$, we find that

$$
\begin{aligned}
& (s-1) \int_{0}^{\infty} \frac{e^{-(s-1) t}}{t+1} d t \\
& \begin{aligned}
\leq(s-1) \int_{0}^{(s-1)^{-1}} \frac{d t}{t+1}+\frac{s-1}{s} & \int_{(s-1)^{-1}}^{\infty}(s-1) e^{-(s-1) t} d t \\
& =(s-1) \log \frac{s}{s-1}+\frac{s-1}{s} e^{-1}
\end{aligned}
\end{aligned}
$$

It follows that as $s \downarrow 1$, the above error term tends to zero, and so $F(s) \rightarrow$ $\gamma-B_{1}$ as desired.

## 5. Primes and probability

In $\S 1$ we discussed how Gauss was led to the prime number theorem by observing that the "local density" of primes near is $x$ is approximately $1 / \log x$. This observation can be used to support many additional statements about primes, the majority of which seem to lie very deep.

We can get a feeling for the reasoning involved in these heuristic arguments by considering a quantitative version of a problem discussed qualitatively in Chapter 1. Suppose that $a \bmod m$ is a (fixed) coprime residue class: How many primes $p \leq x$ are there with $p \equiv a(\bmod m)$ ? Denote the answer to this question by $\pi(x ; m, a)$. In Chapter 1 we mentioned the theorem of Dirichlet that there are always infinitely many such primes, i.e., that $\pi(x ; m, a) \rightarrow \infty$. Now we would like to know how quickly $\pi(x ; m, a)$ tends to infinity.

The numbers not exceeding $x$ from the residue class $a \bmod m$ have the form $a+m r$, where $r \lesssim x / m$. The Gauss philosophy says that a number
chosen at random near $a+m r$ should be prime with probability about $1 / \log (a+m r)$. So, parroting our reasoning in $\S 1$, we might conjecture that $\pi(x ; m, a) \approx \int_{1}^{x / m} d t / \log (a+m t)$. But this cannot be correct: It is easy to check (e.g., using L'Hôpital's rule) that the integral here is asymptotic to $\operatorname{Li}(x) / m$. But there are only $\varphi(m)$ coprime residue classes modulo $m$, so if our guess is correct, then summing over the coprime residue classes modulo $m$ accounts for only $\sim(\varphi(m) / m) \operatorname{Li}(x)$ primes $p \leq x$. Since $\varphi(m) / m<1$ when $m>1$, this contradicts the prime number theorem.

Where did we go wrong? The answer is in our pretending that $a+m r$ is a typical number of its size. Suppose $p$ is prime. Loosely speaking, the probability that a number near $a+m r$ is a multiple of $p$ is $1 / p$. What is the probability that $a+m r$ itself is a multiple of $p$ ? If $p$ does not divide $m$, then the congruence $a+m r \equiv 0(\bmod p)$ has exactly one solution $r$ modulo $p$, and so again this probability is $1 / p$. But if $p$ does divide $m$, then $p$ never divides a number of the form $a+m r$, and so $a+m r$ has a leg up on being prime over its neighbors.

To account for this we introduce a correction factor $c_{p}$ for each prime $p$, defined as a ratio of two probabilities: In the numerator of $c_{p}$ we put the probability that $a+m r$ is not divisible by $p$, and in the denominator we put the probability that a typical number near $a+m r$ is not divisible by $p$. Then $c_{p}=1$ for primes $p$ not dividing $m$, while $c_{p}=(1-1 / p)^{-1}$ when $p$ does divide $m$. Each $c_{p}$ measures the leg up that a number of the form $a+m r$ has over its neighbors, as seen from the perspective of $p$. The Chinese remainder theorem suggests that these effects modulo $p$ should be treated as independent, which in turns suggests that our earlier guesstimate for $\pi(x ; m, a)$ should be multiplied by a factor of $\prod_{p} c_{p}=m / \varphi(m)$. This leads to the new prediction that when $\operatorname{gcd}(a, m)=1$, we have

$$
\pi(x ; m, a) \sim \frac{1}{\varphi(m)} \operatorname{Li}(x) \quad(x \rightarrow \infty)
$$

Unlike our former guess, this is no longer obviously false, and in fact it can be proved correct by the same methods used to establish the prime number theorem. It is known as the prime number theorem for arithmetic progressions.

Let's try something harder: How many $n \leq x$ are there for which both $n$ and $n+2$ are prime? This quantity is traditionally denoted $\pi_{2}(x)$. The Gauss philosophy suggests that a random pair of integers "near $n$ " should be simultaneously prime with probability about $1 /(\log n)^{2}$. But $n$ and $n+2$ do not form a typical pair of integers "near $n$ ". Indeed, let $p$ be a prime number. The probability that neither element of a pair of random numbers near $n$ is divisible by $p$ is $(1-1 / p)^{2}$. But the probability that neither $n$ nor

Table 3. Comparison of $\pi_{2}(x)$ and $L_{2}(x):=2 C_{2} \int_{2}^{x} \frac{d t}{(\log t)^{2}}$. The last column gives the percentage error, computed as $\left|L_{2}(x)-\pi_{2}(x)\right| / \pi_{2}(x)$.

| $x$ | $\pi_{2}(x)$ | $L_{2}(x)-\pi_{2}(x)$ | \% error |
| ---: | ---: | ---: | ---: |
| $10^{5}$ | 1,224 | 25 | $2.0 \%$ |
| $10^{6}$ | 8,169 | 79 | $9.7 \times 10^{-1} \%$ |
| $10^{7}$ | 58,980 | -226 | $3.8 \times 10^{-1} \%$ |
| $10^{8}$ | 440,312 | 56 | $1.3 \times 10^{-2} \%$ |
| $10^{9}$ | $3,424,506$ | 802 | $2.3 \times 10^{-2} \%$ |
| $10^{10}$ | $27,412,679$ | $-1,262$ | $4.6 \times 10^{-3} \%$ |
| $10^{11}$ | $224,376,048$ | $-7,183$ | $3.2 \times 10^{-3} \%$ |
| $10^{12}$ | $1,870,585,220$ | $-25,353$ | $1.4 \times 10^{-3} \%$ |
| $10^{13}$ | $15,834,664,872$ | $-66,567$ | $4.2 \times 10^{-4} \%$ |
| $10^{14}$ | $135,780,321,665$ | $-56,771$ | $4.2 \times 10^{-5} \%$ |
| $10^{15}$ | $1,177,209,242,304$ | $-750,443$ | $6.4 \times 10^{-5 \%}$ |

$n+2$ is divisible by $p$ is $(1-\nu(p) / p)$, where

$$
\nu(p):=\#\{n \bmod p: n(n+2) \equiv 0(\bmod p)\} .
$$

For each prime $p$, put $c_{p}:=(1-\nu(p) / p)(1-1 / p)^{-2}$. Then we might expect that $\pi_{2}(x) \approx\left(\prod_{p} c_{p}\right) \int_{2}^{x} \frac{d t}{(\log t)^{2}}$. Noting that $\nu(p)=1$ if $p=2$ and $\nu(p)=2$ if $p>2$, this conjecture becomes:

Conjecture 3.18 (Twin prime conjecture, quantitative form). As $x \rightarrow \infty$, we have $\pi_{2}(x) \sim 2 C_{2} \int_{2}^{x} \frac{d t}{(\log t)^{2}}$, where $C_{2}:=\prod_{p>2}\left(1-(p-1)^{-2}\right)$.

The constant $C_{2}$ is called the twin prime constant. The numerical evidence for Conjecture 3.18 is very persuasive; see Table 3.

As discussed in Chapter 1, the twin prime conjecture can be viewed as a special case of Schinzel's Hypothesis H. We can now formulate a quantitative version of that general conjecture. Suppose that $f_{1}(T), \ldots, f_{r}(T) \in \mathbf{Z}[T]$ are $r$ distinct polynomials with integer coefficients, that each $f_{i}(T)$ has a positive leading coefficient, that each is irreducible over $\mathbf{Z}$, and that

$$
\begin{equation*}
\text { there is no prime } p \text { dividing } f_{1}(n) \cdots f_{r}(n) \text { for every } n \in \mathbf{Z} \text {. } \tag{3.27}
\end{equation*}
$$

Let $d_{i}$ denote the degree of $f_{i}$. Then $\log \left|f_{i}(n)\right|$ is asymptotic to $d_{i} \log |n|$ as $n \rightarrow \infty$. Our heuristic suggests that

$$
\pi_{f_{1}, \ldots, f_{r}}(x):=\#\left\{n \leq x: f_{1}(n), \ldots, f_{r}(n) \text { are simultaneously prime }\right\}
$$

should be asymptotic to

$$
\begin{equation*}
C\left(f_{1}, \ldots, f_{r}\right) \frac{1}{d_{1} \cdots d_{r}} \int_{2}^{x} \frac{d t}{(\log t)^{r}}, \tag{3.28}
\end{equation*}
$$

where

$$
\begin{equation*}
C\left(f_{1}, \ldots, f_{r}\right):=\prod_{p} \frac{1-\nu(p) / p}{(1-1 / p)^{r}} \tag{3.29}
\end{equation*}
$$

and

$$
\nu(p):=\#\left\{n \bmod p: f_{1}(n) \cdots f_{r}(n) \equiv 0(\bmod p)\right\} .
$$

Notice that the condition (3.27) amounts to the assertion that $\nu(p)<p$ for every prime $p$.

It is worth taking a step back to see if this conjecture makes sense. Does the infinite product (3.29) even converge? This is not at all obvious, even in the simple case when $r=1$. Nevertheless, as shown by Bateman \& Horn [BH62], this is true: The product (3.29) always converges; in fact it always converges to a positive real number. The proof uses some elementary results of Landau on the distribution of prime ideals in algebraic number fields. The positivity of the constant $C\left(f_{1}, \ldots, f_{r}\right)$ means that this quantitative formulation of Hypothesis H really does imply the qualitative formulation of Chapter 1.

The basic argument of this section has many other applications. We close this section with two examples that do not fall under the rubric of Hypothesis H.

For a positive integer $N$, let $R(N)$ be the number of (ordered) pairs of primes $p$ and $q$ for which $p+q=N$. A well-known conjecture of Goldbach asserts that $R(N)>0$ whenever $N>2$ is even. The methods of this section suggest much more:

Conjecture 3.19 (Goldbach conjecture, quantitative form). As $N \rightarrow \infty$ through even numbers, we have

$$
\begin{equation*}
R(N) \sim 2 C_{2}\left(\prod_{p \mid N, p>2} \frac{p-1}{p-2}\right) \int_{2}^{N} \frac{d t}{(\log t)^{2}} . \tag{3.30}
\end{equation*}
$$

Here $C_{2}$ is the twin prime constant.
We leave the task of justifying this conjecture as Exercise 4.
For our last example we consider the distribution of Mersenne primes, i.e., primes of the form $2^{p}-1$. A number near $2^{p}-1$ is prime with probability roughly $1 / \log \left(2^{p}-1\right) \approx 1 /(p \log 2)$. But $2^{p}-1$ is atypical in that we can rule out small prime divisors in advance: If $q$ is a prime divisor of $2^{p}-1$, then 2 has order $p$ modulo $q$, which implies that $q \equiv 1(\bmod p)$. In particular, every prime divisor of $2^{p}-1$ is at least $p$.

Let us make the working assumption that this is the only relevant difference between $2^{p}-1$ and a number typical for its size. Since a typical
integer is divisible by a prime $q$ with probability $1 / q$, this suggests that we multiply our former probability $1 /(p \log 2)$ by $\prod_{q \leq p}(1-1 / q)^{-1}$. By Mertens' theorem, $\prod_{q \leq p}(1-1 / q)^{-1} \sim e^{\gamma} \log p($ as $p \rightarrow \infty)$. This suggests that among the primes $p \leq x$, we should expect

$$
\approx \sum_{p \leq x} e^{\gamma} \frac{\log p}{p \log 2}=\frac{e^{\gamma}}{\log 2} \sum_{p \leq x} \frac{\log p}{p} \sim \frac{e^{\gamma}}{\log 2} \log x
$$

for which $2^{p}-1$ is also prime. (Here we have used (3.25) to estimate the last sum.) So we arrive at the following prediction:
Conjecture 3.20. There are infinitely many primes p for which $2^{p}-1$ is prime. In fact, the number of such $p \leq x$ is asymptotic to $c \log x$ where $c=e^{\gamma} / \log 2$ and $\gamma$ is the Euler-Mascheroni constant.

## Notes

The discussion in $\S 1$ of Gauss's discovery of the prime number theorem is based on [LeV96]. With all due respect to Gauss's ingenuity and industriousness, it must be admitted that Gauss's observations do not provide any explanation for the truth of the prime number theorem. A candidate for such an explanation was proposed by Hawkins [Haw58].

To explain Hawkins's idea, we first recall the classical sieve of Eratosthenes for obtaining a list of the prime numbers: Begin with the sequence $2,3,4,5, \ldots$ of natural numbers $n>1$. Circle the first uncircled number $m$ on the list. Now remove from the list every $n>m$ which is divisible by $m$. If this process is repeated indefinitely, the sequence of circled numbers coincides with the set of primes.

Suppose, following Hawkins, that the deterministic removal step above is replaced with the following random step: Instead of removing each $n>m$ which is divisible by $m$, remove each $n>m$ with probability $1 / m$. That is, for each $n>m$, roll an $m$-sided die (with faces labeled " 1 " thru " $m$ "), and remove the number $n$ if the toss comes up " 1 " and keep the number $n$ otherwise. In this case, indefinite repetition results in a random sequence $\mathcal{P}$. Let $\pi_{\mathcal{P}}(x)$ be the number of terms of $\mathcal{P}$ not exceeding $x$. The following remarkable theorem was conjectured by Hawkins ([Haw74], but see already [Erd65, p. 213]) and proved by Wunderlich [Wun75]:
А Theorem 3.21. With probability 1 , we have $\pi_{\mathcal{P}}(x) \sim x / \log x$ as $x \rightarrow \infty$.
Informally, this result says that Eratosthenes-like sieves tend to produce sequences which satisfy the conclusion of the prime number theorem - so maybe it should not come as a shock that the sequence actually produced by the sieve of Eratosthenes has this property. A story with a similar moral is told in [GLMU56, HB58].

Table 4. Comparison of $\pi(x)$ and $E(x):=\operatorname{Li}(x)-\pi(x)$ along powers of 10 , from $x=10^{14}$ through $x=10^{23} . E(x)$ is shown rounded to the nearest integer.

| $x$ | $\pi(x)$ | $E(x)$ |
| ---: | ---: | ---: |
| $10^{14}$ | $3,204,941,750,802$ | 314,891 |
| $10^{15}$ | $29,844,570,422,669$ | $1,052,617$ |
| $10^{16}$ | $279,238,341,033,925$ | $3,214,631$ |
| $10^{17}$ | $2,623,557,157,654,233$ | $7,956,588$ |
| $10^{18}$ | $24,739,954,287,740,860$ | $21,949,554$ |
| $10^{19}$ | $234,057,667,276,344,607$ | $99,877,774$ |
| $10^{20}$ | $2,220,819,602,560,918,840$ | $222,744,643$ |
| $10^{21}$ | $21,127,269,486,018,731,928$ | $597,394,253$ |
| $10^{22}$ | $201,467,286,689,315,906,290$ | $1,932,355,207$ |
| $10^{23}$ | $1,925,320,391,606,803,968,923$ | $7,250,186,215$ |

If we take a careful look at Table 2, we are led to wonder whether the prime number theorem is not too modest an assertion. Put

$$
E(x):=\operatorname{Li}(x)-\pi(x)
$$

The prime number theorem asserts that $E(x)=o(\pi(x))$, while the data in Table 2 suggests that $E(x)$ is actually of a much smaller order of magnitude than $x$. In Table 4 we extend the comparison of $\pi(x)$ and $\operatorname{Li}(x)$ up to $10^{23}$. Inspecting this table, we find that the numbers in the third column are only about half the length of those in the second, which suggests that perhaps $|E(x)| \lesssim \sqrt{\pi(x)}$. While nothing of this sort can yet be proved, this behavior is not unexpected: It has been known since Riemann that the size of $E(x)$ is intimately connected with the location of the zeros of $\zeta(s)$. The so-called Riemann Hypothesis asserts that all the nonreal zeros of $\zeta(s)$ lie on the line $\Re(s)=1 / 2$. As shown by von Koch [Koc01] in 1901, the Riemann Hypothesis is equivalent to the bound

$$
E(x)=O(\sqrt{x} \log x)
$$

Unfortunately, we still cannot even prove that $E(x)=O\left(x^{1-\epsilon}\right)$ for a fixed positive value of $\epsilon$. The best-known result is (in somewhat rough form) that for each fixed $\alpha<3 / 5$, there is a constant $C_{\alpha}>0$ with

$$
\begin{equation*}
E(x) \ll x \exp \left(-C_{\alpha}(\log x)^{\alpha}\right) \tag{3.31}
\end{equation*}
$$

That this is the state-of-the-art reflects an embarrassing lack of twentieth century progress, since the result (3.31) with $\alpha=1 / 2$ was established by de la Vallée-Poussin [VP99] already in 1899.

In the opposite direction, it is known that von Koch's conditional bound on $E(x)$, if correct, is close to best possible:
$\star$ Theorem 3.22 (Littlewood [Lit14]). There are constants $c^{-}<0<c^{+}$ for which the following holds: There is a sequence of $x$ tending to infinity along which

$$
E(x)>c^{+} x^{1 / 2} \log \log \log x / \log x
$$

and a sequence of $x$ tending to infinity along which

$$
E(x)<c^{-} x^{1 / 2} \log \log \log x / \log x .
$$

Littlewood's theorem is usually quoted in connection with one of its more surprising consequences, namely that $E(x)$ changes sign infinitely often. (Tables 2 and 4 might lead one to the contrary conjecture that $E(x) \rightarrow \infty$ as $x \rightarrow \infty$.)

Our proofs of the theorems of Chebyshev and Mertens incorporate a number of later simplifications. For a discussion of these authors' original methods, one should consult the beautiful monograph of Narkiewicz [Nar04], in particular, Chapter 3. This monograph is the source of much of the historical content throughout this book.

The quantitative forms of the twin prime and Goldbach conjectures which we discussed in $\S 5$ are due to Hardy \& Littlewood [HL23]. Their approach was considerably more complicated than ours; the realization that conjectures of this type could be derived from the "Gauss philosophy" on the local density of primes appears to be due to Selmer [Sel42] (see also [Gol60]). Bateman \& Horn [BH62] were the first to suggest, in full generality, the quantitative form of Hypothesis H discussed in §5. Conjecture 3.20 was suggested independently by Pomerance, Selfridge and Wagstaff (see, e.g., [Wag83]).

## Exercises

1. Let $A$ be a set of natural numbers and let $A(x):=\#\{a \leq x: a \in A\}$. Show that if $\sum_{a \in A} a^{-1}$ converges, then $A$ has asymptotic density zero.
2. (a) (Golomb [Gol62]) Show that for each integer $k>1$, there is at least one natural number $n$ for which $n / \pi(n)=k$.
(b) Show that the set of $n$ for which $\pi(n)$ divides $n$ has asymptotic density zero. (Cf. [EP90].)
3. Should one expect that there are infinitely many primes of the form $n!+1$ ? What about $p!+1$, where $p$ itself is prime?
4. Provide a convincing argument suggesting the truth of Conjecture 3.19.
5. Using only the divergence of $\sum_{p} p^{-1}$, show that $\lim \sup _{x \rightarrow \infty} \frac{\pi(x)}{x /(\log x)^{1+\epsilon}}$ is infinite for each fixed $\epsilon>0$.
6. (a) Suppose $\left\{a_{n}\right\}_{n \geq 1}$ and $\left\{b_{n}\right\}_{n \geq 1}$ are sequences of real numbers where $a_{n} \rightarrow \infty$ and $a_{n} \sim b_{n}$ as $n \rightarrow \infty$. Show that $a_{n} \log a_{n} \sim b_{n} \log b_{n}$ as $n \rightarrow \infty$.
(b) Write $p_{n}$ for the $n$th prime number. Taking $a_{n}:=p_{n} / \log p_{n}$ and $b_{n}:=n$, deduce from the prime number theorem that $p_{n} \sim n \log n$ as $n \rightarrow \infty$.
7. (Continuation) Prove that $p_{n+1} / p_{n} \rightarrow 1$ as $n \rightarrow \infty$. Show also that $\{p / q: p, q$ prime $\}$ is a dense subset of $(0, \infty)$.
8. Show that if $m$ is a fixed natural number, then $\operatorname{Li}(x)$ may be estimated as

$$
\frac{x}{\log x}+\frac{x}{(\log x)^{2}}+\frac{2 x}{(\log x)^{3}}+\cdots+\frac{(m-1)!x}{(\log x)^{m}}+O_{m}\left(\frac{x}{(\log x)^{m+1}}\right) .
$$

Assuming (3.31), show that the same expansion is valid for $\pi(x)$ replacing $\operatorname{Li}(x)$.
9. (Landau [Lan01]) Let $\pi^{\prime}(x)$ be the number of primes in the interval $(x, 2 x]$. Assuming the prime number theorem, show that $\pi^{\prime}(x) \sim \pi(x)$ as $x \rightarrow \infty$. Assuming (3.31), show that $\pi(x)>\pi^{\prime}(x)$ for large $x$, and that in fact $\pi(x)-\pi^{\prime}(x) \rightarrow \infty$ as $x \rightarrow \infty$.

Remark. It is tempting to conjecture, as Hardy \& Littlewood did in 1923 (see [HL23]), that the interval ( $0, x$ ] always contains at least as many primes as the interval $(y, x+y]$ whenever $x, y \geq 2$. However, this is probably false; Hensley \& Richards [HR73] have shown that it contradicts the prime $k$-tuples conjecture, which is a special case of Schinzel's Hypothesis H.
10. (Gelfond \& Schnirelmann [Gel46]; cf. [Mon94, Chapter 10]) Show that for each natural number $N$,

$$
\operatorname{lcm}[1,2, \ldots, N]=\exp (\psi(N))
$$

Deduce that the expression

$$
e^{\psi(2 N+1)} \int_{0}^{1} x^{N}(1-x)^{N} d x
$$

represents a positive integer, and use this to give another proof that $\psi(x) \geq x \log 2+O(\log x)$ as $x \rightarrow \infty$.
11. (Brun [Bru17]) For $x \geq 2$, let $N=N(x)$ be the number of natural numbers $n \leq x$ divisible by some prime $p \in(\sqrt{x}, x]$.
(a) Noting that each natural number $n \leq x$ can be divisible by at most one prime $p \in(\sqrt{x}, x]$, show that $N \geq \sum_{\sqrt{x}<p \leq x}\lfloor x / p\rfloor$.
(b) Deduce from the trivial bound $N \leq x$ that $\sum_{\sqrt{x}<p \leq x} 1 / p \leq 2$.
(c) Use the result of (b) to give another proof that $\sum_{p \leq x} p^{-1} \ll \log \log x$ as $x \rightarrow \infty$.
12. In this exercise and the next we establish Bertrand's postulate in its full strength: For every positive integer $n$, there is a prime $p$ with $n<$ $p \leq 2 n$. The proof described here is a hybrid of Ramanujan's argument (described in §3.3) and an argument of Erdős [Erd32], and can be found in [Sha83, §9.3C].
(a) Check that $\prod_{n+1<p \leq 2 n+1} \left\lvert\,\binom{ 2 n+1}{n+1}\right.$ for every integer $n \geq 0$.
(b) Prove that $\binom{2 n+1}{n+1} \leq 4^{n}$ for each integer $n \geq 0$.
(c) Use (a) and (b) to fashion an inductive proof that $\prod_{p \leq N} p \leq 4^{N}$ for all nonnegative integers $N$. Thus $\theta(x) \leq 2 x \log 2$ for all $x \geq 0$.
(d) Check that $\binom{2 n+1}{n+1}$ is divisible by every prime $p \leq n+1$ which possesses a power belonging to the interval ( $n+1,2 n+1]$. Use this to show that $\exp (\psi(N)) \leq 4^{N}$ for every natural number $N \geq 0$. Thus $\psi(x) \leq 2 x \log 2$ for every $x \geq 0$.

Remark. The argument of (a)-(c) is due to Erdős \& Kalmár (see [Erd89]). Erdős's 1932 paper had a more complicated proof of a slightly weaker bound for $\theta(x)$.
13. (Continuation) Recall that for each $x \geq 0$, we have

$$
\begin{equation*}
T(x)-2 T(x / 2) \leq \psi(x)-\psi(x / 2)+\psi(x / 3) . \tag{3.32}
\end{equation*}
$$

(a) Show that if $n$ is a nonnegative integer, then $\binom{2 n}{n} \geq 4^{n} /(2 n+1)$. Hint: What does the $2 n$th row of Pascal's triangle look like?
(b) Show that

$$
\sum_{\substack{n<p^{k} \leq 2 n \\ k \geq 2}} \log p \leq \sqrt{2 n} \log \sqrt{2 n}
$$

(c) Deduce from (3.32) (with $x=2 n$ ) and (d) of the last exercise that

$$
\theta(2 n)-\theta(n) \geq \frac{1}{3} n \log 4-\log (2 n+1)-\sqrt{2 n} \log \sqrt{2 n}
$$

(d) Conclude from (c) that there is always a prime in the interval ( $n, 2 n$ ] whenever $n \geq 82$.
(e) The primes $2,3,5,7,13,23,43,83$ form a sequence with each less than twice the next. Use this to argue that there is always a prime in the interval $(n, 2 n]$ for $n<82$ as well.
14. (Richert [Ric49]) Using the full form of Bertrand's postulate, show that every integer $n>6$ can be written as a sum of distinct prime numbers. Hint: Start by observing that if $6<n \leq 19$, then $n$ is a sum of distinct primes $\leq 11$.
15. Let $p_{1}=2, p_{2}=3, p_{3}=5, \ldots$ be the sequence of primes and put $d_{n}:=p_{n+1}-p_{n}$. Deduce from Theorem 3.5 that $\lim \inf d_{n} / \log p_{n}<\infty$ and $\limsup d_{n} / \log p_{n}>0$.

Remark. The twin prime conjecture says that $d_{n}=2$ infinitely often, which of course implies that

$$
\begin{equation*}
\liminf _{n \rightarrow \infty} d_{n} / \log p_{n}=0 \tag{3.33}
\end{equation*}
$$

In 2005, Goldston, Pintz, and Yıldırım (see [GPY, GMPY06] and the survey [Sou07]) proved that (3.33) holds unconditionally, which had been a long-standing open problem. Their method can be elaborated on to show that on an infinite set of $n$,

$$
d_{n} \ll\left(\log p_{n}\right)^{1 / 2}\left(\log \log p_{n}\right)^{2} .
$$

The principal tool needed in their argument is a theorem of Bombieri and Vinogradov. Roughly speaking, the Bombieri-Vinogradov theorem asserts that the primes are as well-distributed in arithmetic progressions, on average, as the Generalized Riemann Hypothesis predicts for each individual progression. A plausible strengthening of the BombieriVinogradov conjecture, due to Elliott \& Halberstam, would imply that infinitely often $d_{n} \leq 16$, which would put us agonizingly close to the twin prime conjecture. In fact, any improvement of the Bombieri-Vinogradov theorem in the direction of the Elliott-Halberstam conjecture would imply the existence of a constant $C$ with $d_{n} \leq C$ infinitely often. However, such improvements seem to lie very deep.

In the opposite direction, it was shown by Westzynthius [Wes31] already in 1931 that $\lim \sup _{n \rightarrow \infty} \frac{d_{n}}{\log p_{n}}=\infty$. The best result in this direction is due to Erdős [Erd35b] and Rankin [Ran38]: For some constant $c>0$ and infinitely many $n$,

$$
d_{n}>c \log p_{n} \frac{\log \log p_{n} \log \log \log \log p_{n}}{\left(\log \log \log p_{n}\right)^{2}}
$$

According to work of Pintz [Pin97], we can take $c=2 e^{\gamma}$. Erdős offered a prize of $\$ 10,000$ for a proof that $c$ could be taken arbitrarily large.
16. (Continuation; Erdős \& Turán [ET48])
(a) Prove that $d_{n}<d_{n+1}$ for infinitely many $n$.
(b) Prove that $d_{n}>d_{n+1}$ for infinitely many $n$. Hint: Assume that $d_{n} \leq d_{n+1}$ whenever $n \geq N_{0}$. Fix $C>0$ so that $d_{m}<C \log p_{m}$ for infinitely many $m$. Show that there is a $k_{0} \in \mathbf{N}$ with the property that if $k$ is a natural number with $k \geq k_{0}$, then $d_{n}=k$ can hold for at most $k$ consecutive values of $n$. Now argue that if $d_{m}<C \log p_{m}$, then $p_{m+1}-2=\sum_{i=1}^{m} d_{i} \ll\left(\log p_{m}\right)^{3}$.
Remark. Open problems about $d_{n}$ abound; here are two: Is $d_{n}=d_{n+1}$ for infinitely many $n$ ? Is $d_{n}<d_{n+1}<d_{n+2}$ infinitely often?
17. Show that the series $\sum_{n=1}^{\infty} \frac{1}{n\left(p_{n+1}-p_{n}\right)^{\lambda}}$ diverges when $\lambda=1$, and give a heuristic argument suggesting that it diverges for every real $\lambda$.
18. For each integer $n>1$, let $P(n)$ denote the largest prime factor of $n$. Determine the set of real numbers $\lambda$ for which $\sum_{n>1} \frac{1}{n^{\lambda} P(n)}$ converges.
19. (Sierpiński [Sie64]) It is an easy consequence of Hypothesis H that for every positive integer $k$, there are infinitely many primes of the form $n^{2}+k$. Show (unconditionally) that for every natural number $N$, there is a positive integer $k$ for which there are at least $N$ primes of the form $n^{2}+k$. Hint: For every $p$, one can write $p=\lfloor\sqrt{p}\rfloor^{2}+k$ for some $k \ll \sqrt{p}$.
20. Show that for every $N \in \mathbf{N}$, there is an even integer $k>0$ for which there are at least $N$ prime pairs $p, p+k$.
21. (Mertens, Lindqvist \& Peetre [LP97]) In this exercise we derive an alternative expression for the constant $B_{1}$ in Theorem 3.14, namely

$$
\begin{equation*}
B_{1}=\gamma+\sum_{n=2}^{\infty} \frac{\mu(n)}{n} \log \zeta(n) \tag{3.34}
\end{equation*}
$$

(Using the expansion (3.34) and a table of $\zeta$-values compiled by Legendre, Mertens showed that $B_{1}=0.2614972128 \ldots$..) By the results of §4, in order to prove (3.34) it is enough to show that

$$
\begin{equation*}
B_{2}=-\sum_{n=2}^{\infty} \frac{\mu(n)}{n} \log \zeta(n) \tag{3.35}
\end{equation*}
$$

(Here $B_{2}=\sum_{k \geq 2} \sum_{p} \frac{1}{k p^{k}}$, as in the statement of Theorem 3.15.)
(a) Prove that for real $s>1$, we have $s^{-1} \log \zeta(s)=\int_{2}^{\infty} \frac{\pi(t)}{t\left(t^{s}-1\right)} d t$.

Show also that $B_{2}=\int_{2}^{\infty} \frac{\pi(t)}{t^{2}(t-1)} d t$.
(b) Prove that for $|x|<1$,

$$
\frac{x^{2}}{1-x}=-\sum_{m=2}^{\infty} \mu(m) \frac{x^{m}}{1-x^{m}} .
$$

(c) Taking $x=1 / t$ in part (b), deduce that for $t>1$,

$$
\frac{1}{t^{2}(t-1)}=-\sum_{m=2}^{\infty} \frac{\mu(m)}{t\left(t^{m}-1\right)}
$$

(d) Use the results of (a)-(c) to prove (3.35).
22. (Pomerance [Pom79]) Using $p_{n}$ to denote the $n$th prime number, let $G$ be the collection of points $\left(n, p_{n}\right) \in \mathbf{R}^{2}$, where $n \in \mathbf{N}$. We call $G$ the prime number graph.
(a) Show that every line in $\mathbf{R}^{2}$ contains only finitely many points of $G$.
(b) In the remainder of this exercise we prove that there are lines in the plane which contain arbitrarily many points of $G$. For this we may replace $G$ by $G^{\prime}:=\left\{\left(p_{n}, n\right): n \in \mathbf{N}\right\}$.
Let $k \in \mathbf{N}$. Put $u=e^{k}, v=u+u / \log u$, and let $T$ be the parallelogram bounded by the vertical lines $x=u, x=v$ and the diagonal lines with slope $1 / k$ through $\left(u, \operatorname{Li}(u)+2 u /(\log u)^{4}\right)$ and $\left(u, \operatorname{Li}(u)-3 u /(\log u)^{4}\right)$. Prove that there are $\ll k u /(\log u)^{4}$ lines of slope $1 / k$ passing through lattice points contained in $T$ (as $k \rightarrow \infty)$.
(c) Assuming that $\pi(x)-\operatorname{Li}(x)=o\left(x /(\log x)^{4}\right)$ (which follows from (3.31)), prove that every point $\left(p_{n}, n\right)$ with $u \leq p_{n} \leq v$ lies inside $T$ once $k$ is sufficiently large.
(d) Show that as $k \rightarrow \infty$, there are $\gg u /(\log u)^{2}$ points $\left(p_{n}, n\right)$ with $u \leq p_{n} \leq v$. Conclude from (b) and (c) that there is a line of slope $1 / k$ passing through $\gg \frac{1}{k}(\log u)^{2}=k$ of these points.
23. (Hardy \& Ramanujan [HR17], Turán [Tur34]) Write $\omega(n)$ for the number of distinct prime factors of $n$ and $\Omega(n)$ for the number of prime factors of $n$ counted with multiplicity. (Thus, if $n=\prod_{i=1}^{k} p_{i}^{e_{i}}$, where the $p_{i}$ are distinct primes and each $e_{i} \geq 1$, then $\omega(n)=k$ and $\Omega(n)=\sum_{i=1}^{k} e_{i}$.)
(a) Show that for $x \geq 3$, we have $\sum_{n \leq x} \omega(n)=x \log \log x+O(x)$ and $\sum_{n \leq x} \omega(n)^{2}=x(\log \log x)^{2}+O(x \log \log x)$.
(b) Deduce from (a) that $\sum_{n \leq x}(\omega(n)-\log \log x)^{2}=O(x \log \log x)$.
(c) Conclude from (b) that if $B>0$, then the number of $n \leq x$ with $|\omega(n)-\log \log x|>B \sqrt{\log \log x}$ is $\ll x / B^{2}$, where the implied
constant is absolute. Hence $\omega(n)$ is very close to $\log \log x$ for most $n \leq x$.
(d) Show that $\sum_{n \leq x}(\Omega(n)-\omega(n))^{2}=O(x)$, and deduce that the result of (c) holds with $\omega$ replaced by $\Omega$.

Remark. For fixed real numbers $B_{1}<B_{2}$, a beautiful theorem of Erdős \& Kac [EK40] asserts that

$$
\frac{1}{x} \#\left\{n \leq x: B_{1} \leq \frac{\omega(n)-\log \log x}{\sqrt{\log \log x}} \leq B_{2}\right\} \rightarrow \frac{1}{\sqrt{2 \pi}} \int_{B_{1}}^{B_{2}} e^{-u^{2} / 2} d u
$$

as $x \rightarrow \infty$, and the same with $\omega$ replaced by $\Omega$. Actually the ErdősKac result is far more general and can be viewed as an analogue of the central limit theorem for additive arithmetic functions. ${ }^{1}$ The Erdős-Kac theorem stands with the Erdős-Wintner theorem (discussed in the notes to Chapter 8) as one of the foundational results in probabilistic number theory.
24. (Continuation; Erdős [Erd55, Erd60]) Suppose $N$ is a natural number. The $N \times N$ multiplication table is defined as the $N \times N$ array whose $i$ th row, $j$ th column entry is $i \cdot j$. Since multiplication is commutative, it is clear that the number $A(N)$ of distinct entries in this table is bounded by the number of unordered pairs of integers from $[1, N]$, which is just $\frac{1}{2} N(N+1)$. The following rough argument suggests that $A(N)$ is considerably smaller:

For most ordered pairs of integers $(i, j)$ with $1 \leq i, j \leq N$, the number $\Omega(i \cdot j)=\Omega(i)+\Omega(j)$ of prime factors of $i \cdot j$ is very close to $2 \log \log N$ by Exercise 23. But most numbers $n \leq N^{2}$ have about $\log \log \left(N^{2}\right) \sim \log \log N$ prime factors. So the multiplication table contains mostly atypical numbers, and so it cannot contain very many of the numbers $n \leq N^{2}$.
Fill in the details of this argument to construct a rigorous proof that $A(N) / N^{2} \rightarrow 0$ as $N \rightarrow \infty$.
Remark. As a consequence of a detailed study of the distribution of divisors of natural numbers, Ford [For08a] (see also [For08b]) proved that

$$
A(N) \asymp \frac{N^{2}}{(\log N)^{\delta}(\log \log N)^{3 / 2}}, \quad \text { where } \quad \delta:=1-\frac{1+\log \log 2}{\log 2} .
$$

25. (Erdős [Erd79]) Define $\omega(n ; z):=\sum_{p \mid n ; p \leq z} 1$, so that $\omega(n)=\omega(n ; n)$.
(a) Show that if $x \geq z \geq 3$, then

$$
\sum_{n \leq x}(\omega(n ; z)-\log \log z)^{2} \ll x \log \log z
$$

[^2](b) Define a sequence of positive real numbers $\left\{z_{j}\right\}_{j=1}^{\infty}$ by putting $z_{j}:=$ $\exp \left(\exp \left(j^{4}\right)\right)$. Show that if $x \geq z_{j}$, then there are $\ll x j^{-2}$ natural numbers $n \leq x$ with $\left|\omega\left(n ; z_{j}\right)-\log \log z_{j}\right|>\left(\log \log z_{j}\right)^{3 / 4}$.
(c) Now let $\epsilon>0$. Show that one can choose a positive real number $Z$, depending only on $\epsilon$, so that the following holds: If $x$ is sufficiently large, then all but at most $\epsilon x$ natural numbers $n \leq x$ satisfy $|\omega(n ; z)-\log \log z|<40(\log \log z)^{3 / 4}$ for all $Z<z \leq x$.
(d) Prove that all of the assertions of (a)-(c) remain valid if $\omega(n ; z)$ is replaced by $\Omega(n ; z):=\sum_{p^{k} \mid n, p \leq z} 1$.
26. (Continuation) For $n \in \mathbf{N}$ and $1 \leq k \leq \omega(n)$, let $p_{k}(n)$ denote the $k$ th smallest prime divisor of $n$. Show that for each $\epsilon>0$ and $\eta>0$, there is a natural number $K$ for which the following holds: The set of natural numbers $n$ for which
$$
k(1-\epsilon)<\log \log p_{k}(n)<k(1+\epsilon)
$$
for every $K<k \leq \omega(n)$ has lower density at least $1-\eta$. Roughly speaking, this says that for large $k$, the $k$ th prime factor of a typical natural number is approximately $e^{e^{k}}$.
27. The twin prime conjecture illustrates how difficult it can be to control the multiplicative structure of neighboring integers. In this exercise we give an elementary example where this is possible.
(a) Define a sequence of finite subsets $S_{i} \subset \mathbf{N}$ as follows: Let $S_{2}=$ $\{2,3\}$. Assuming $S_{r}$ has already been defined, let $M$ be the product of all the elements of $S_{r}$ and put $S_{r+1}:=\{M\} \cup\left\{M-a: a \in S_{r}\right\}$. Check that for each $r$, the set $S_{r}$ has $r$ elements and $\left|a_{1}-a_{2}\right|=$ $\operatorname{gcd}\left(a_{1}, a_{2}\right)$ for every pair of distinct elements $a_{1}, a_{2} \in S_{r}$. (This important construction is due to Heath-Brown [HB87].)
(b) Suppose that $f: \mathbf{N} \rightarrow \mathbf{C}^{\times}$is a completely multiplicative arithmetic function and that its image $f(\mathbf{N})$ is finite. Show that the set of $n \in \mathbf{N}$ for which $f(n)=f(n+1)$ has positive lower density. Hint: Choose a natural number $r>|f(\mathbf{N})|$, and list the elements $a_{1}<$ $a_{2}<\cdots<a_{r}$ of $S_{r}$. Put $M=\prod_{i=1}^{r} a_{i}$. Start by observing that for any $k \in \mathbf{N}$, at least two of the values $\left\{f\left(k M+a_{j}\right)\right\}_{1 \leq j \leq r}$ must coincide.
(c) Using (b), show that for each fixed $m \in \mathbf{N}$, a positive proportion of natural numbers $n$ satisfy $\Omega(n) \equiv \Omega(n+1)(\bmod m)$.

Remark. For further results on the multiplicative structure of consecutive integers, see Hildebrand's elegant survey [Hil97].
28. (Montgomery \& Wagon [MW06]) Suppose that $W(x)$ is a real-valued function of $x$ which is decreasing for $x \geq 2$. Prove that if

$$
\int_{2}^{x} W(t) \log t \frac{d t}{t} \sim \log x
$$

then $W(x) \sim 1 / \log x$ as $x \rightarrow \infty$. Hint: Obtain a lower bound for $\liminf _{x \rightarrow \infty} W(x) \log x$ by observing that

$$
W(x) \int_{x}^{x^{1+\epsilon}} \log t \frac{d t}{t} \geq \int_{x}^{x^{1+\epsilon}} W(t) \log t \frac{d t}{t} \sim \epsilon \log x .
$$

Replacing the limits of integration with $x^{1-\epsilon}$ and $x$, establish an analogous upper bound for $\lim \sup _{x \rightarrow \infty} W(x) \log x$.
29. (Continuation) We now prove that if $\pi(x) \sim x / L(x)$ for a function $L(x)$ which is positive-valued and increasing for $x \geq 2$, then necessarily $L(x) \sim \log x$, so that the prime number theorem holds. Note that this generalizes Theorem 3.4.

Put $f(x)=x^{-1} \log x$, so that $\sum_{p \leq x} f(p) \sim \log x$ by (3.25).
(a) Show that $\sum_{p \leq x} f(p) \sim-\int_{2}^{x} \pi(t) f^{\prime}(t) d t$ as $x \rightarrow \infty$.
(b) Prove that $\int_{2}^{x} \pi(t) f^{\prime}(t) d t \sim \int_{2}^{x}(t / L(t)) f^{\prime}(t) d t$.
(c) Deduce from (a), (b), and (3.25) that

$$
\int_{2}^{x} L(t)^{-1} \log t \frac{d t}{t} \sim \log x .
$$

(d) Conclude from Exercise 28 that $1 / L(x) \sim 1 / \log x$, so that $L(x) \sim$ $\log x$.

Remark. See Exercise 7.3 for a different strengthening of Theorem 3.4.
30. In this exercise and the next we explore what can be proved with our present tools about the magnitude of the divisor function $\tau(n)$.
(a) Show that $\sum_{n \leq x} \tau(n)=x \log x+O(x)$ for $x \geq 1$. So on average, a natural number $n \leq x$ has about $\log x$ divisors.
(b) Show that $2^{\omega(n)} \leq \tau(n) \leq 2^{\Omega(n)}$ for every natural number $n$. Deduce from Exercise 23 that for each $B>0$, all but $O\left(x / B^{2}\right)$ of the natural numbers $n \leq x$ satisfy

$$
2^{\log \log x-B \sqrt{\log \log x}} \leq \tau(n) \leq 2^{\log \log x+B \sqrt{\log \log x}} .
$$

Since $2^{\log \log x}=(\log x)^{\log 2}$, this shows that most $n \leq x$ have significantly fewer divisors than the average.
31. (Continuation; Wigert [Wig07]) Let $n$ be a natural number not exceed$\operatorname{ing} x$. Let $A:=\prod_{p^{e} \| n, p \leq \frac{\log x}{(\log \log x)^{2}}} p^{e}$ and put $B:=\prod_{p^{e} \| n, p>\frac{\log x}{(\log \log x)^{2}}} p^{e}$.
(a) Show that $\tau(A) \leq 2^{O\left(\log x /(\log \log x)^{2}\right)}$ as $x \rightarrow \infty$.
(b) Show that $\Omega(B) \leq(1+o(1)) \log x / \log \log x$. Deduce that $\tau(B) \leq$ $2^{(1+o(1)) \frac{\log x}{\log \log x}}$.
(c) Conclude from (a) and (b) that $\tau(n) \leq 2^{(1+o(1)) \frac{\log x}{\log \log x}}$.
(d) By considering the product of an initial segment of the primes, show that there is a sequence of $n$ tending to infinity along which

$$
\tau(n) \geq 2^{(1+o(1)) \frac{\log n}{\log \log n}} .
$$

Thus (c) is best possible. You may assume the prime number theorem for this part of the exercise, but this is not necessary.
32. Recall that $\Psi(x, y)$ denotes the number of $y$-smooth $n \leq x$, i.e., the number of natural numbers $n \leq x$ all of whose prime divisors are $\leq y$. Rankin [Ran38] observed that for any $\sigma>0$, one has

$$
\begin{equation*}
\Psi(x, y) \leq \sum_{\substack{n \leq x \\ p \mid n \Rightarrow p \leq y}}\left(\frac{x}{n}\right)^{\sigma}=x^{\sigma} \prod_{p \leq y}\left(1-p^{-\sigma}\right)^{-1} \tag{3.36}
\end{equation*}
$$

Suppose now that $x \geq y \geq 2$, and put $\sigma:=1-\frac{1}{2 \log y}$. Show that

$$
\frac{1}{p^{\sigma}}-\frac{1}{p} \ll \frac{\log p}{p \log y}
$$

uniformly for primes $p \leq y$, and deduce that the product appearing in (3.36) is $\ll \log y$. Conclude that for $x \geq y \geq 2$,

$$
\Psi(x, y) \ll x e^{-u / 2} \log y, \quad \text { where } \quad u:=\frac{\log x}{\log y}
$$

and the implied constant is absolute.
33. (Gauss's polynomial prime number theorem) For each $A(T) \in \mathbf{F}_{q}[T]$, put $|A|:=q^{\operatorname{deg} A}$. Define the zeta function $\zeta_{q}(s)$ of $\mathbf{F}_{q}[T]$ by setting $\zeta_{q}(s):=\sum_{A}|A|^{-s}$, where $A$ runs over all monic polynomials in $\mathbf{F}_{q}[T]$. Let $\pi(q ; n)$ denote the number of monic irreducible polynomials of degree $n$ over $\mathbf{F}_{q}$.
(a) Show that for $s>1$, we have $\zeta_{q}(s)=1 /\left(1-q^{1-s}\right)$.
(b) Show that for $s>1$, there is a product representation of $\zeta_{q}(s)$, namely $\zeta_{q}(s)=\prod_{P}\left(1-|P|^{-s}\right)^{-1}$, where $P$ runs over all monic irreducible polynomials in $\mathbf{F}_{q}[T]$.
(c) From (a) and (b), deduce that with $u=q^{-s}$,

$$
\frac{1}{1-q u}=\prod_{j=1}^{\infty}\left(\frac{1}{1-u^{j}}\right)^{\pi(q ; j)} .
$$

(d) Starting with the result of (c), show that

$$
\begin{equation*}
\sum_{d \geq 1} d \pi(q ; d) \frac{u^{d}}{1-u^{d}}=\frac{q u}{1-q u} . \tag{3.37}
\end{equation*}
$$

Hint: Take the logarithmic derivative.
(e) By comparing the coefficients of $u^{n}$ on both sides of (3.37), deduce that $q^{n}=\sum_{d \mid n} d \pi(q ; d)$. Conclude that $\pi(q ; n)=\frac{1}{n} \sum_{d \mid n} \mu(d) q^{n / d}$.
(f) Show that $\left|\pi(q ; n)-q^{n} / n\right| \leq 2 q^{n / 2} / n$ for every prime power $q$ and every natural number $n$.
If we set $X=q^{n}$, then we have just shown that $\pi(q ; n)$ is very close to $X / \log _{q} X$, where $\log _{q}$ denotes the logarithm with base $q$. This is strikingly reminiscent of the prime number theorem.
34. (Mertens' theorem for polynomials) Show that $\prod_{\operatorname{deg} P \leq n}(1-1 /|P|)=$ $e^{-\gamma} / n+O(1 / n)$, where $\gamma$ is the Euler-Mascheroni constant. Here $n$ is a natural number, $P$ runs over the monic irreducible polynomials in $\mathbf{F}_{q}[T]$ of degree at most $n$, and the implied constant is understood to be absolute (independent of both $q$ and $n$ ). Proceed as follows:
(a) Reduce the proof to the assertion that

$$
\sum_{\operatorname{deg} P \leq n} \sum_{k \geq 1} \frac{1}{k|P|^{k}}=\log n+\gamma+O(1 / n)
$$

(b) Use the results of Exercise 33 to show that we have the (exact) identity

$$
\sum_{P, k: \operatorname{deg} P^{k} \leq n} \frac{1}{k|P|^{k}}=\sum_{m \leq n} \frac{1}{m} .
$$

(c) Complete the proof by first estimating $\sum_{m \leq n} m^{-1}$ using Lemma 3.16 and then showing that

$$
\sum_{\operatorname{deg} P \leq n} \sum_{k>n / \operatorname{deg} P} \frac{1}{k|P|^{k}} \ll \frac{1}{n} .
$$

This argument is due to K. Conrad (see [EHM02]).
35. (A polynomial analogue of the twin prime conjecture) Capelli's theorem (proved, e.g., as [Lan02, Theorem 9.1]) asserts that if $F$ is an arbitrary field, $a \in F$ and $n \in \mathbf{N}$, then the binomial $T^{n}-a$ is irreducible in $F[T]$ unless one of the following holds:
(i) there is a prime $l$ dividing $n$ for which $a$ is an $l$ th power in $F$,
(ii) 4 divides $n$ and $a=-4 b^{4}$ for some $b \in F$.

Using this result, show that $T^{3^{k}}-3$ and $T^{3^{k}}-2$ are both irreducible over $\mathbf{F}_{7}$ for every integer $k \geq 0$. In particular, there are infinitely many monic polynomials $A(T) \in \mathbf{F}_{7}[T]$ for which $A$ and $A+1$ are both irreducible.

If you are feeling ambitious, prove that this last claim holds with $\mathbf{F}_{7}$ replaced by any finite field with more than 3 elements. This result is due to Hall [Hal03, Hal06].

Remark. Actually this claim holds even for the field $\mathbf{F}_{3}$, but a somewhat different argument is required. For this and other generalizations, see [Pol08a]. See also [Eff08], [Pol08b].

## Primes in Arithmetic Progressions

> When Gauss says he has proved something, it is very probable ... when Cauchy says it, you can bet equally well pro or contra, but when Dirichlet says it, it is certain. I prefer to leave myself out of this Delikatessen. C. G. J. Jacobi, letter to von Humboldt

## 1. Introduction

In this chapter we prove Dirichlet's result [Dir37, Dir39, Dir41] that if $a$ and $m$ are integers with $m>0$ and $\operatorname{gcd}(a, m)=1$, then there are infinitely many primes $p \equiv a(\bmod m)$. Actually, we shall prove more, namely that for $x \geq 4$,

$$
\begin{equation*}
\sum_{\substack{p \leq x \\=a(\bmod m)}} \frac{\log p}{p}=\frac{1}{\varphi(m)} \log x+O(1) \tag{4.1}
\end{equation*}
$$

where the implied constant may depend on $m$. The infinitude of primes $p \equiv a(\bmod m)$ is of course an obvious consequence, but (4.1) says much more. In light of (3.25), we can view (4.1) as an equidistribution statement, asserting that (in a peculiar average sense) the fraction of primes falling into a given coprime residue class is exactly $1 / \varphi(m)$. Moreover, as shown in Exercise 2, the estimate (4.1) implies that

$$
\begin{equation*}
\pi(x ; m, a) \gg_{a, m} \frac{x}{\log x}, \tag{4.2}
\end{equation*}
$$

which can be considered an analogue of Chebyshev's lower bound on $\pi(x)$ from Chapter 3.

As an application of Dirichlet's result, we close the chapter with a proof of Legendre's characterization of the integers expressible as a sum of three squares.

## 2. Progressions modulo 4

We begin by considering the case when $m=4$. Define a function $\chi: \mathbf{Z} \rightarrow \mathbf{C}$ by putting

$$
\chi(n):= \begin{cases}(-1)^{(n-1) / 2} & \text { if } 2 \nmid n  \tag{4.3}\\ 0 & \text { otherwise } .\end{cases}
$$

It is straightforward to check that $\chi(a b)=\chi(a) \chi(b)$ for every pair of integers $a, b$. So, at least formally (i.e., ignoring issues of convergence),

$$
\begin{equation*}
\prod_{p}\left(1-\frac{\chi(p)}{p}\right)^{-1}=\sum_{n \geq 1} \frac{\chi(n)}{n} \tag{4.4}
\end{equation*}
$$

(cf. Theorem 1.2). Let $L$ denote the right-hand series; then

$$
\begin{align*}
L: & =1-\frac{1}{3}+\frac{1}{5}-\frac{1}{7}+\frac{1}{9}-\frac{1}{11}+\cdots  \tag{4.5}\\
& =(1-1 / 3)+(1 / 5-1 / 7)+(1 / 9-1 / 11)+\cdots>2 / 3 .
\end{align*}
$$

In particular, $L>0$. Taking the logarithm of both sides of (4.4), we deduce that as $x \rightarrow \infty$,

$$
\sum_{\substack{p^{k} \leq x \\ p^{k} \equiv 1(\bmod 4)}} \frac{1}{k p^{k}}-\sum_{\substack{p^{k} \leq x \\ p^{k} \equiv 3(\bmod 4)}} \frac{1}{k p^{k}}=\log L+o(1) .
$$

The terms corresponding to $k \geq 2$ contribute a negligible amount to both sums, which implies that

$$
\sum_{\substack{p \leq x \\ p \equiv 3(\bmod 4)}} \frac{1}{p}-\sum_{\substack{p \leq x \\ p \equiv 1(\bmod 4)}} \frac{1}{p}
$$

is $O(1)$. Since $\sum_{p \leq x} p^{-1} \sim \log \log x$ (by Mertens' first theorem), both $\sum_{p \leq x, p \equiv 1(\bmod 4)} p^{-1}$ and $\sum_{p \leq x, p \equiv 3(\bmod 4)} p^{-1}$ are $\sim \frac{1}{2} \log \log x$. In particular, both coprime residue classes modulo 4 contain infinitely many primes.

Unfortunately, it is by no means apparent how to justify the identity (4.4). (Our only tool for establishing a factorization like (4.4) is Theorem 1.2 , but its hypotheses do not hold in this case.) There are various ways to work around this; the most common is to replace the series $\sum \chi(n) n^{-1}$ with $\sum \chi(n) n^{-s}$, where $s>1$. Then $\sum \chi(n) n^{-s}$ is absolutely convergent,
and so from Theorem 1.2 we obtain the analogue of (4.4). Following the above argument, we now find that $\sum_{p} \chi(p) p^{-s}$ remains bounded as $s \downarrow 1$. Since $\sum_{p} p^{-s}$ diverges to infinity as $s \downarrow 1$, it must be that 1 and -1 both occur as the value of $\chi(p)$ for infinitely many primes $p$. This again shows that both coprime progressions modulo 4 contain infinitely many primes. We will follow a different route in this text; rather than alter the terms of the series $\sum \chi(n) n^{-1}$, we alter the range of summation, working with the truncations $\sum_{n \leq x} \chi(n) n^{-1}$.

Suppose now that $m$ is any natural number and $a \in \mathbf{Z}$. Then

$$
\begin{aligned}
\sum_{\substack{n \leq x \\
n \equiv a(\bmod m)}} \frac{\Lambda(n)}{n} & =\sum_{\substack{p^{k} \leq x \\
p^{k} \equiv a(\bmod m)}} \frac{\log p}{p^{k}} \\
& =\sum_{\substack{p \leq x \\
p \equiv a(\bmod m)}} \frac{\log p}{p}+\sum_{k \geq 2} \sum_{\substack{p \leq x^{1 / k} \\
p^{k} \equiv a(\bmod m)}} \frac{\log p}{p^{k}} .
\end{aligned}
$$

By (3.24), the double sum here is absolutely bounded. Consequently,

$$
\begin{equation*}
\sum_{\substack{p \leq x \\ p \equiv a(\bmod m)}} \frac{\log p}{p}=\sum_{\substack{n \leq x \\ n \equiv a(\bmod m)}} \frac{\Lambda(n)}{n}+O(1) . \tag{4.6}
\end{equation*}
$$

Thus estimates for $\sum \log p / p$, taken over the primes $p \equiv a(\bmod m)$, follow from estimates for $\sum \Lambda(n) / n$, taken over natural numbers $n \equiv a(\bmod m)$.

Now specialize again to the case $m=4$. Let $\chi$ be as defined in (4.3), and let $\chi_{0}$ be the indicator function of the odd integers. Then $\chi_{0}+\chi$ is twice the characteristic function of the arithmetic progression $1 \bmod 4$, and $\chi_{0}-\chi$ is twice the characteristic function of the arithmetic progression $3 \bmod 4$. This suggests studying the summatory functions

$$
\begin{equation*}
\sum_{n \leq x} \frac{\chi_{0}(n) \Lambda(n)}{n} \quad \text { and } \quad \sum_{n \leq x} \frac{\chi(n) \Lambda(n)}{n} \tag{4.7}
\end{equation*}
$$

The first of these behaves very much like the sum $\sum_{n \leq x} \Lambda(n) / n$ investigated in Chapter 3:

$$
\begin{align*}
\sum_{n \leq x} \frac{\chi_{0}(n) \Lambda(n)}{n} & =\sum_{n \leq x} \frac{\Lambda(n)}{n}-\sum_{2^{k} \leq x} \frac{\log 2}{2^{k}} \\
& =\sum_{n \leq x} \frac{\Lambda(n)}{n}+O(1)=\log x+O(1) \tag{4.8}
\end{align*}
$$

the final equality coming from (3.22). To understand the second sum appearing in (4.7), we notice that $L$, defined in (4.5), is an alternating series
with terms decreasing in absolute value. Thus, if we use $N$ to denote the smallest odd integer exceeding $x$, then for every $x \geq 1$,

$$
\begin{equation*}
\left|\sum_{n>x} \frac{\chi(n)}{n}\right| \leq\left|\frac{\chi(N)}{N}\right|=\frac{1}{N}<\frac{1}{x} \tag{4.9}
\end{equation*}
$$

Following Mertens, we observe next that

$$
\begin{aligned}
\sum_{n \leq x} \frac{\chi(n) \log n}{n} & =\sum_{n \leq x} \frac{\chi(n)}{n} \sum_{d \mid n} \Lambda(d) \\
& =\sum_{d \leq x} \Lambda(d) \sum_{\substack{n \leq x \\
d \mid n}} \frac{\chi(n)}{n} \\
& =\sum_{d e \leq x} \frac{\chi(d e) \Lambda(d)}{d e}=\sum_{d \leq x} \frac{\chi(d) \Lambda(d)}{d} \sum_{e \leq x / d} \frac{\chi(e)}{e}
\end{aligned}
$$

The inner sum here is equal to $L-\sum_{e>x / d} \chi(e) e^{-1}=L+O(d / x)$. Substituting this above tells us that

$$
\begin{aligned}
\sum_{n \leq x} \frac{\chi(n) \log n}{n} & =L \sum_{d \leq x} \frac{\chi(d) \Lambda(d)}{d}+O\left(\frac{1}{x} \sum_{d \leq x} \Lambda(d)\right) \\
& =L \sum_{d \leq x} \frac{\chi(d) \Lambda(d)}{d}+O(1)
\end{aligned}
$$

since $\sum_{d \leq x} \Lambda(d)=\psi(x) \ll x$. Also, $\sum_{n \leq x} \chi(n) \log n / n=O(1)$, since

$$
\frac{\log 1}{1}-\frac{\log 3}{3}+\frac{\log 5}{5}-\cdots
$$

is an alternating series with eventually decreasing terms. Thus

$$
L \sum_{d \leq x} \frac{\chi(d) \Lambda(d)}{d}=O(1)
$$

and since $L \neq 0$, it follows that

$$
\begin{equation*}
\sum_{d \leq x} \frac{\chi(d) \Lambda(d)}{d}=O(1) \tag{4.10}
\end{equation*}
$$

From (4.8) and (4.10), we deduce that

$$
\begin{aligned}
& \sum_{\substack{n \leq x \\
n \equiv 1(\bmod 4)}} \frac{\Lambda(n)}{n}+\sum_{\substack{n \leq x \\
n \equiv 3(\bmod 4)}} \frac{\Lambda(n)}{n}=\log x+O(1), \\
& \sum_{\substack{n \leq x \\
n \equiv 1(\bmod 4)}} \frac{\Lambda(n)}{n}-\sum_{\substack{n \leq x \\
n \equiv 3(\bmod 4)}} \frac{\Lambda(n)}{n}=O(1) .
\end{aligned}
$$

Adding these estimates shows that

$$
\sum_{\substack{n \leq x \\ n \equiv 1(\bmod 4)}} \frac{\Lambda(n)}{n}=\frac{1}{2} \log x+O(1),
$$

and subtracting yields the same result for $n$ restricted to the residue class $3 \bmod 4$. Referring to equation (4.6) shows that the same estimates hold for the sums $\sum \log p / p$. This completes the proof of (4.1) when $m=4$.

In general, to prove Dirichlet's theorem for all coprime progressions modulo $m$, we will need to consider $\varphi(m)-1$ series analogous to the single series $L$ appearing in this proof. The most difficult part of the argument consists of showing that none of these series converges to zero.

Remark. For the remainder of this chapter, up until the exercises, we adopt the convention that all implied constants (unless otherwise stated) may depend on $m$. Further dependence will be mentioned explicitly.

## 3. The characters of a finite abelian group

To carry out the strategy which proved successful for progressions modulo 4 , we first need to understand the appropriate analogues of the function $\chi$, as defined in (4.3), for a general modulus $m$. These turn out to be the Dirichlet characters modulo $m$, which arise in a natural way from the characters of the unit group $(\mathbf{Z} / m \mathbf{Z})^{\times}$.
3.1. The classification of characters. Let $G$ be a finite abelian group (written multiplicatively). By a character of $G$ we mean a homomorphism

$$
\chi: G \rightarrow \mathbf{C}^{\times},
$$

i.e., a function from $G$ to the nonzero complex numbers satisfying

$$
\begin{equation*}
\chi(a b)=\chi(a) \chi(b) \tag{4.11}
\end{equation*}
$$

for every $a, b \in G$. The set of characters of $G$ is denoted $\hat{G}$. We let $\chi_{0}$ denote the trivial character which is identically 1 . Note that if $\chi$ is a character of $G$, then every value which $\chi$ assumes is a root of unity. Indeed, if the order
of $g \in G$ is $n$, then $\chi(g)^{n}=\chi\left(g^{n}\right)=\chi(1)=1$, so that $\chi(g)$ is an $n$th root of unity.

Our goal in this section is to classify the characters of an arbitrary finite abelian group $G$. We first treat the case when $G$ is cyclic. Fix a generator $g_{0}$ of $G$. The value of $\chi\left(g_{0}\right)$ determines $\chi(g)$ for every $g \in G$; indeed, if $g=g_{0}^{k}$, then $\chi(g)=\chi\left(g_{0}^{k}\right)=\chi\left(g_{0}\right)^{k}$. From the preceding paragraph, $\chi\left(g_{0}\right)$ must be a $|G|$ th root of unity, and so $G$ has at most $|G|$ characters. Moreover, we see that there are precisely $|G|$ characters if and only if for every $|G|$ th root of unity $\eta$, there is a character $\chi$ of $G$ with $\chi\left(g_{0}\right)=\eta$. And it is easy to describe a character $\chi$ of $G$ for which this holds: Simply define $\chi$ by putting $\chi\left(g_{0}^{k}\right)=\eta^{k}$ for all $k$. $\chi$ is well-defined, since if $g=g_{0}^{k_{1}}=g_{0}^{k_{2}}$, then $k_{1} \equiv k_{2}(\bmod |G|)$, so that $\eta^{k_{1}}=\eta^{k_{2}}$. Moreover, it is straightforward to verify (4.11) in this case, so that $\chi$ is a genuine character of $G$. We have thus achieved a complete classification of the characters of a finite cyclic group.

An arbitrary finite abelian group of course need not be cyclic, but according to a well-known classification theorem, every such group is a direct sum of cyclic groups. In other words, one can always find elements $g_{1}, \ldots, g_{k} \in G$ with respective orders $n_{1}, \ldots, n_{k}$ (say), with the property that every $g \in G$ has a unique representation in the form

$$
g_{1}^{e_{1}} g_{2}^{e_{2}} \cdots g_{k}^{e_{k}}, \quad \text { where } 0 \leq e_{i}<n_{i} \text { for each } 1 \leq i \leq k .
$$

If $\chi$ is a character of $G$, then $\chi$ is completely determined by $\chi\left(g_{1}\right), \ldots, \chi\left(g_{k}\right)$. Since $\chi\left(g_{i}\right)$ must be an $n_{i}$ th root of unity for each $i$, we see that there are at most $\prod_{i=1}^{k} n_{i}=|G|$ characters of $G$. Moreover, if for each $1 \leq i \leq k$ we let $\eta_{i}$ be an arbitrary $n_{i}$ th root of unity, then it is easy to check that putting

$$
\begin{equation*}
\chi\left(g_{1}^{e_{1}} \cdots g_{k}^{e_{k}}\right):=\eta_{1}^{e_{1}} \eta_{2}^{e_{2}} \cdots \eta_{k}^{e_{k}} \tag{4.12}
\end{equation*}
$$

gives us a well-defined character $\chi$ of $G$ with $\chi\left(g_{i}\right)=\eta_{i}$ for each $1 \leq i \leq k$. So there are precisely $|G|$ elements of $\hat{G}$, and we understand them all.

Remark. For the purposes of this chapter, we need not invoke any classification results from group theory. We only need to understand the case when $G=(\mathbf{Z} / m \mathbf{Z})^{\times}$. In this case the existence of a decomposition into cyclic groups is elementary: Indeed, the Chinese remainder theorem guarantees that if $m=\prod_{i=1}^{k} p_{i}^{e_{i}}$, then $(\mathbf{Z} / m \mathbf{Z})^{\times} \cong \prod_{i=1}^{k}\left(\mathbf{Z} / p_{i}^{e_{i}} \mathbf{Z}\right)^{\times}$, and we obtain the desired decomposition of $G$ once we recall that (see, e.g., [IR90, Theorems $\left.2,2^{\prime}\right]$ )

$$
\left(\mathbf{Z} / p^{e} \mathbf{Z}\right)^{\times} \cong \begin{cases}\mathbf{Z} / 2 \mathbf{Z} \oplus \mathbf{Z} / 2^{e-2} \mathbf{Z} & \text { if } p=2, e>2 \\ \mathbf{Z} /\left((p-1) p^{e-1}\right) \mathbf{Z} & \text { otherwise }\end{cases}
$$

3.2. The orthogonality relations. The characters of a finite abelian group obey certain orthogonality relations, which play an essential role in the proof of Dirichlet's theorem. In the situation that concerns us, when $G=(\mathbf{Z} / m \mathbf{Z})^{\times}$, these relations allow us to express the characteristic function of a coprime residue class modulo $m$ as a linear combination of characters.

Before stating these relations, we note that $\hat{G}$ can be made into a group (called the dual group of $G$ ) by defining, for $\chi, \psi \in \hat{G}$,

$$
(\chi \psi)(g):=\chi(g) \psi(g),
$$

i.e., by defining the multiplication pointwise. The trivial character $\chi_{0}$ now serves as the identity. Associativity and commutativity follow from the corresponding properties of $\mathbf{C}^{\times}$. And inverses are easy; for each $\chi \in \hat{G}$, define $\chi^{-1}$ by putting

$$
\chi^{-1}(g):=\chi(g)^{-1} .
$$

The right-hand side exists since $\chi$ takes values in the nonzero complex numbers, and the homomorphism property of $\chi^{-1}$ follows from inverting both sides of (4.11). Notice that because the values $\chi$ assumes are always roots of unity, we have $\chi^{-1}=\bar{\chi}$, where $\bar{\chi}$ is defined by $\bar{\chi}(g):=\overline{\chi(g)}$ for each $g \in G$.

Now suppose that $\chi \in \hat{G}$ is nontrivial, i.e., $\chi \neq \chi_{0}$. Then there is an element $h \in G$ with $\chi(h) \neq 1$. Since $G$ is a group, $h g$ runs over the elements of $G$ as $g$ does. Thus, setting $S_{\chi}=\sum_{g \in G} \chi(g)$, one has

$$
\chi(h) S_{\chi}=\chi(h) \sum_{g \in G} \chi(g)=\sum_{g \in G} \chi(h g)=\sum_{g \in G} \chi(g)=S_{\chi} .
$$

Since $\chi(h) \neq 1$, we must have $S_{\chi}=0$. Thus

$$
\sum_{g \in G} \chi(g)= \begin{cases}|G| & \text { if } \chi=\chi_{0}  \tag{4.13}\\ 0 & \text { otherwise }\end{cases}
$$

Since $\bar{\chi}=\chi^{-1}$ for any character $\chi$, this can be recast as follows: If $\chi$ and $\psi$ are two characters of $G$, then

$$
\sum_{g \in G} \bar{\chi}(g) \psi(g)= \begin{cases}|G| & \text { if } \chi=\psi  \tag{4.14}\\ 0 & \text { otherwise }\end{cases}
$$

Equation (4.14) is the first of two orthogonality relations for characters. It was obtained by studying $\sum \chi(g)$, where $\chi \in \hat{G}$ is fixed and $g$ runs over the elements of the group $G$. To obtain the second orthogonality relation, we investigate the same sum where instead $g \in G$ is fixed and $\chi$ runs over the elements of the group $\hat{G}$. To proceed we require the following lemma:

Lemma 4.1. Let $G$ be a finite abelian group and let $g \neq 1$ be an element of $G$. Then there exists a character $\chi \in \hat{G}$ with $\chi(g) \neq 1$.

Proof. Let $g_{1}, \ldots, g_{k}$ be a system of independent generators for $G$ as in §3.1, so that every element of $G$ admits a unique representation in the form (4.12). Since $g$ is not the identity of $G$, in its representation in the form (4.12) there is at least one exponent $e_{i}$ with $0<e_{i}<n_{i}$. Fix such an $i$, and let $\chi$ be the character of $G$ defined by $\chi\left(g_{1}^{e_{1}} \cdots g_{k}^{e_{k}}\right)=\eta_{i}^{e_{i}}$, where $\eta_{i}$ is a fixed primitive $n_{i}$ th root of unity. Then $\chi(g) \neq 1$.

Now let $g \neq 1$ be an element of $G$ and choose $\psi \in \hat{G}$ with $\psi(g) \neq 1$. Set $S_{g}=\sum_{\chi \in \hat{G}} \chi(g)$. Since $\hat{G}$ forms a group, $\psi \chi$ runs over all elements of $\hat{G}$ as $\chi$ does. Consequently,

$$
\psi(g) S_{g}=\psi(g) \sum_{\chi \in \hat{G}} \chi(g)=\sum_{\chi \in \hat{G}}(\psi \chi)(g)=\sum_{\chi \in \hat{G}} \chi(g)=S_{g} .
$$

Hence

$$
\sum_{\chi \in \hat{G}} \chi(g)= \begin{cases}|G| & \text { if } g=1 \\ 0 & \text { otherwise }\end{cases}
$$

Noting that for each $g \in G$,

$$
\chi\left(g^{-1}\right)=\chi(g)^{-1}=\overline{\chi(g)}=\bar{\chi}(g),
$$

we find that

$$
\sum_{\chi \in \hat{G}} \bar{\chi}(g) \chi(h)= \begin{cases}|G| & \text { if } g=h,  \tag{4.15}\\ 0 & \text { otherwise } .\end{cases}
$$

This is the second orthogonality relation.
3.3. Dirichlet characters. Let $m$ be a natural number and let $G=$ $(\mathbf{Z} / m \mathbf{Z})^{\times}$, the group of units modulo $m$. For each $\chi \in \hat{G}$, we introduce an associated function $\tilde{\chi}$ defined on the set of integers coprime to $m$ by putting

$$
\tilde{\chi}(a):=\chi(a \bmod m) .
$$

We extend $\tilde{\chi}$ to a function defined on all of $\mathbf{Z}$ by setting $\tilde{\chi}(a):=0$ whenever $\operatorname{gcd}(a, m)>1$. The functions $\tilde{\chi}$ are known as the Dirichlet characters modulo $m$. Instead of continuing to write " $\tilde{\chi}$ ", in what follows we adopt a customary abuse of notation and use the same symbol $\chi$ for both the function on $G$ and the associated function on $\mathbf{Z}$.

It is easy to see that every Dirichlet character $\chi$ modulo $m$ has both of the following properties:
(i) $\chi$ is periodic modulo $m$, i.e., $\chi(a+m)=\chi(a)$ for every $a \in \mathbf{Z}$.
(ii) $\chi$ is completely multiplicative, i.e., for every $a, b \in \mathbf{Z}$,

$$
\chi(a b)=\chi(a) \chi(b) .
$$

Moreover, the Dirichlet characters obey the following orthogonality relations:

Theorem 4.2. Let $m$ be a positive integer and let $\chi$ and $\psi$ be two Dirichlet characters modulo $m$. Then

$$
\sum_{a \bmod m} \bar{\chi}(a) \psi(a)= \begin{cases}\varphi(m) & \text { if } \chi=\psi^{-1}  \tag{4.16}\\ 0 & \text { otherwise } .\end{cases}
$$

Theorem 4.3. Let $m$ be a positive integer. If $a, b \in \mathbf{Z}$ and $\operatorname{gcd}(a, m)=1$, then

$$
\sum_{\chi} \bar{\chi}(a) \chi(b)= \begin{cases}\varphi(m) & \text { if } a \equiv b(\bmod m)  \tag{4.17}\\ 0 & \text { otherwise } .\end{cases}
$$

Here the sum is over all Dirichlet characters $\chi$ modulo $m$.
These results follow from (4.14) and (4.15) if $G$ is taken to be the $\varphi(m)$ element group $(\mathbf{Z} / m \mathbf{Z})^{\times}$: Theorem 4.2 is immediate from (4.14), since the values of $a$ with $\operatorname{gcd}(a, m)=1$ do not contribute to the left-hand side of (4.16). To prove Theorem 4.3, notice that (4.17) follows immediately from (4.15) in the case when $\operatorname{gcd}(a, m)=\operatorname{gcd}(b, m)=1$. If, however, $\operatorname{gcd}(b, m)>1$, then the left-hand side of (4.17) vanishes because $\chi(b)=0$. Since $\operatorname{gcd}(b, m)>1$ implies that $a \not \equiv b(\bmod m)$, the theorem holds in this case as well. (This is where we need the condition in Theorem 4.3 that $\operatorname{gcd}(a, m)=1$.)

## 4. The $L$-series at $s=1$

To each Dirichlet character $\chi$ we associate the Dirichlet $L$-series

$$
\begin{equation*}
L(s, \chi):=\sum_{n=1}^{\infty} \frac{\chi(n)}{n^{s}} . \tag{4.18}
\end{equation*}
$$

For our purposes, only the series corresponding to nontrivial characters are of interest and these are only of interest at $s=1$. Nevertheless, because there is no extra difficulty involved, we begin by treating the series corresponding to nontrivial Dirichlet characters whenever $s>0$.

Lemma 4.4. Let $\chi$ be a nontrivial Dirichlet character modulo $m$. Then (4.18) converges for every $s>0$. Moreover, for every $s>0$ and $x \geq 1$,

$$
\left|\sum_{n>x} \frac{\chi(n)}{n^{s}}\right| \leq 2 \varphi(m) x^{-s} .
$$

In particular, $\sum_{n>x} \chi(n) n^{-1} \ll x^{-1}$.

Proof. Put $S(x)=\sum_{n \leq x} \chi(n)$. Theorem 4.2 implies that $\sum \chi(n)$ vanishes when taken over any block of $m$ consecutive integers, which in turn shows that $|S(x)| \leq \varphi(m)$ for every $x$. By partial summation,

$$
\begin{equation*}
\sum_{n \leq x} \frac{\chi(n)}{n^{s}}=\frac{S(x)}{x^{s}}+\int_{1}^{x} s \frac{S(t)}{t^{s+1}} d t \tag{4.19}
\end{equation*}
$$

As $x \rightarrow \infty$, the first term on the right goes to 0 , since $S(x)$ remains bounded while $x^{s}$ tends to infinity. The last factor converges as $x \rightarrow \infty$, by comparison with the absolutely convergent integral $\int_{1}^{\infty} s \frac{\varphi(m)}{t^{s+1}} d t=\varphi(m)$. This proves the first claim.

To bound the tail of $L(s, \chi)$, we apply partial summation once again:

$$
\begin{aligned}
\sum_{n>x} \frac{\chi(n)}{n^{s}} & \left.=\left(\frac{S(y)}{y^{s}}-\frac{S(x)}{x^{s}}+\int_{x}^{y} s \frac{S(t)}{t^{s+1}} d t\right)\right]_{y=\infty} \\
& =-\frac{S(x)}{x^{s}}+\int_{x}^{\infty} s \frac{S(t)}{t^{s+1}} d t
\end{aligned}
$$

The first term is bounded in absolute value by $\varphi(m) x^{-s}$ and the second by $\int_{x}^{\infty} s \frac{\varphi(m)}{t^{s+1}} d t=\varphi(m) x^{-s}$. The stated estimate now follows from the triangle inequality.

## 5. Nonvanishing of $L(1, \chi)$ for complex $\chi$

We say that the Dirichlet character $\chi$ is real if $\chi(\mathbf{Z}) \subset \mathbf{R}$, i.e., if $\chi$ assumes only real values. (In this case, $\chi(\mathbf{Z}) \subset\{0,1,-1\}$, since every nonvanishing value of $\chi$ is a root of unity.) Otherwise, we call $\chi$ a complex character. Our goal in this section is to show that $L(1, \chi)$ is nonvanishing for each complex Dirichlet character $\chi$.

We first connect the vanishing or nonvanishing of $L(1, \chi)$ to the behavior of the partial sums of $\sum \chi(n) \Lambda(n) n^{-1}$.

Lemma 4.5. Let $\chi$ be any nontrivial Dirichlet character modulo m. For $x \geq 4$,

$$
\sum_{n \leq x} \frac{\chi(n) \Lambda(n)}{n}=\left\{\begin{aligned}
O(1) & \text { if } L(1, \chi) \neq 0 \\
-\log x+O(1) & \text { otherwise }
\end{aligned}\right.
$$

Proof when $L(1, \chi) \neq 0$. We mimic the argument of $\S 2$, which corresponds to the case when $\chi$ is the the nontrivial Dirichlet character modulo 4 . We
start by writing

$$
\begin{aligned}
\sum_{n \leq x} \frac{\chi(n) \log n}{n} & =\sum_{n \leq x} \frac{\chi(n)}{n} \sum_{d \mid n} \Lambda(d) \\
& =\sum_{d e \leq x} \frac{\chi(d e) \Lambda(d)}{d e}=\sum_{d \leq x} \frac{\chi(d) \Lambda(d)}{d} \sum_{e \leq x / d} \frac{\chi(e)}{e} .
\end{aligned}
$$

From Lemma 4.4, the inner sum is $L(1, \chi)-\sum_{e>x / d} \chi(e) / e=L(1, \chi)+$ $O(d / x)$. Inserting this above shows that

$$
\begin{aligned}
\sum_{n \leq x} \frac{\chi(n) \log n}{n} & =L(1, \chi) \sum_{d \leq x} \frac{\chi(d) \Lambda(d)}{d}+O\left(\frac{1}{x} \sum_{d \leq x} \Lambda(d)\right) \\
& =L(1, \chi) \sum_{d \leq x} \frac{\chi(d) \Lambda(d)}{d}+O(1)
\end{aligned}
$$

since $\sum_{d \leq x} \Lambda(d)=\psi(x) \ll x$ by (3.17). But we also have

$$
\begin{equation*}
\sum_{n \leq x} \frac{\chi(n) \log n}{n}=O(1) . \tag{4.21}
\end{equation*}
$$

Indeed, with $S(x):=\sum_{n \leq x} \chi(n)$,

$$
\sum_{n \leq x} \frac{\chi(n) \log n}{n}=\frac{S(x) \log x}{x}-\int_{1}^{x} S(t) \frac{1-\log t}{t^{2}} d t
$$

so that (noting that $t^{-1} \log t$ is decreasing for $t \geq e$ )

$$
\left|\sum_{n \leq x} \frac{\chi(n) \log n}{n}\right| \leq \varphi(m) \frac{\log 4}{4}+\varphi(m) \int_{1}^{\infty} \frac{d t}{t^{2}}+\varphi(m) \int_{1}^{\infty} \frac{\log t}{t^{2}} d t \ll 1
$$

Together, (4.20) and (4.21) imply that

$$
L(1, \chi) \sum_{d \leq x} \frac{\chi(d) \Lambda(d)}{d}=O(1) .
$$

Since $L(1, \chi) \neq 0$, the sum here must be bounded (independently of $x$ ), which is the statement of the lemma in this case.

Proof when $L(1, \chi)=0$. Applying Möbius inversion to the relation $\log n=$ $\sum_{d \mid n} \Lambda(d)$, we obtain

$$
\begin{aligned}
\Lambda(n)=\sum_{d \mid n} \mu(d) \log \frac{n}{d} & =\sum_{d \mid n} \mu(d) \log n-\sum_{d \mid n} \mu(d) \log d \\
& =\log n \sum_{d \mid n} \mu(d)-\sum_{d \mid n} \mu(d) \log d=-\sum_{d \mid n} \mu(d) \log d,
\end{aligned}
$$

since for every positive integer $n$, either $\log n=0$ or $\sum_{d \mid n} \mu(d)=0$. So for every $x>0$,

$$
\begin{aligned}
\sum_{d \mid n} \mu(d) \log \frac{x}{d} & =\log x \sum_{d \mid n} \mu(d)+\Lambda(n) \\
& =\left\{\begin{aligned}
\log x+\Lambda(n) & \text { if } n=1 \\
\Lambda(n) & \text { otherwise. }
\end{aligned}\right.
\end{aligned}
$$

Consequently,

$$
\begin{align*}
& \log x+\sum_{n \leq x} \frac{\chi(n) \Lambda(n)}{n}=\sum_{n \leq x} \frac{\chi(n)}{n} \sum_{d \mid n} \mu(d) \log \frac{x}{d} \\
& \quad=\sum_{d \leq x} \mu(d) \log \frac{x}{d} \sum_{\substack{n \leq x \\
d \mid n}} \frac{\chi(n)}{n}=\sum_{d \leq x} \mu(d) \log \frac{x}{d} \frac{\chi(d)}{d} \sum_{e \leq x / d} \frac{\chi(e)}{e} \\
& .22) \quad=L(1, \chi) \sum_{d \leq x} \mu(d)\left(\log \frac{x}{d}\right) \frac{\chi(d)}{d}+R(x), \tag{4.22}
\end{align*}
$$

where (using the estimate of Lemma 4.4)

$$
\begin{align*}
R(x) & \ll \sum_{d \leq x}\left(\log \frac{x}{d}\right) \frac{1}{d} \frac{d}{x}=\frac{1}{x} \sum_{d \leq x}(\log x-\log d) \\
& =\frac{1}{x}(\lfloor x\rfloor \log x-\log \lfloor x\rfloor!) \ll 1 . \tag{4.23}
\end{align*}
$$

(Here we have used Lemma 3.10 to estimate $\log \lfloor x\rfloor$ !.) Since $L(1, \chi)=0$, (4.22) implies that

$$
\log x+\sum_{n \leq x} \frac{\chi(n) \Lambda(n)}{n}=O(1),
$$

which is the assertion of Lemma 4.5 in this case.
We also require an estimate for $\sum_{n \leq x} \chi(n) \Lambda(n) n^{-1}$ when $\chi=\chi_{0}$.
Lemma 4.6. Let $\chi_{0}$ be the trivial character modulo $m$. Then for $x \geq 4$,

$$
\sum_{n \leq x} \frac{\chi_{0}(n) \Lambda(n)}{n}=\log x+O(1)
$$

Proof. Observe that

$$
\sum_{n \leq x} \frac{\Lambda(n)}{n}-\sum_{n \leq x} \frac{\chi_{0}(n) \Lambda(n)}{n}=\sum_{\substack{p \mid m}} \sum_{\substack{p^{k} \leq x \\ k \geq 1}} \frac{\log p}{p^{k}} \leq \sum_{p \mid m} \frac{\log p}{p-1} \ll 1 .
$$

The result now follows from (3.22).

We can now prove the main result of this section.
Theorem 4.7. Let $\chi$ be a complex character modulo $m$. Then $L(1, \chi) \neq 0$.
Proof. Lemmas 4.5 and 4.6 together imply that for $x \geq 4$,

$$
\begin{equation*}
\sum_{\chi} \sum_{n \leq x} \frac{\chi(n) \Lambda(n)}{n}=(1-V) \log x+O(1) \tag{4.24}
\end{equation*}
$$

where $V$ denotes the number of nontrivial $\chi$ with $L(1, \chi)=0$, and the sum is taken over all Dirichlet characters $\chi$ modulo $m$. On the other hand, taking $a=1$ in the orthogonality relation (4.17) shows that

$$
\begin{equation*}
\frac{1}{\varphi(m)} \sum_{\chi} \sum_{n \leq x} \frac{\chi(n) \Lambda(n)}{n}=\sum_{\substack{n \leq x \\ n \equiv 1(\bmod m)}} \frac{\Lambda(n)}{n} \geq 0 . \tag{4.25}
\end{equation*}
$$

If $V>1$, then (4.24) and (4.25) contradict each other for large enough $x$. Thus $V \leq 1$, i.e., $L(1, \chi)$ vanishes for at most one nontrivial character $\chi$.

But if $L(1, \chi)=0$ for some complex character $\chi$, then

$$
0=\overline{L(1, \chi)}=\overline{\sum_{n=1}^{\infty} \frac{\chi(n)}{n}}=\sum_{n=1}^{\infty} \frac{\bar{\chi}(n)}{n}=L(1, \bar{\chi})
$$

also. Since $\chi$ is complex, $\chi \neq \bar{\chi}$, so that $V \geq 2$, a contradiction.
Remarks. For the purpose of demystification, it is worth pointing out that versions of the sums considered in Lemmas 4.5 and 4.6 arise naturally in the analytic context. Indeed,

$$
\begin{aligned}
L(s, \chi)=\prod_{p} & \left(1-\frac{\chi(p)}{p^{s}}\right)^{-1} \quad \text { implies, by logarithmic differentiation, } \\
& \left.-\frac{L^{\prime}(s, \chi)}{L(s, \chi)}=\sum_{n=1}^{\infty} \frac{\chi(n) \Lambda(n)}{n^{s}} \quad \text { (always assuming } \Re(s)>1\right) .
\end{aligned}
$$

The statements of Lemmas 4.5 and 4.6 are also not unexpected: Assume it has been shown that $\zeta(s)$ and $L(s, \chi)$ admit analytic extensions to $\Re(s)>$ 0 , except for simple poles at $s=1$ in the cases of $\zeta(s)$ and the functions $L\left(s, \chi_{0}\right)$. This is a usual first step in the analytic arguments.

If $L(s, \chi)$ is analytic and nonzero at $s=1$, then $\frac{L^{\prime}}{L}(s, \chi)$ is analytic at $s=1$. Suppose, on the other hand, that $L(s, \chi)$ has a zero or pole at $s=1$ (the latter occurring only when $\chi$ is trivial). Let $K$ denote the integer for which $(s-1)^{K} L(s, \chi)$ is analytic and nonzero at $s=1$. Then
$-\frac{L^{\prime}}{L}(s, \chi) \sim \frac{K}{s-1}$ as $s \rightarrow 1$. For $s$ real,

$$
\left|\sum_{n=1}^{\infty} \frac{\chi(n) \Lambda(n)}{n^{s}}\right| \leq \sum_{n=1}^{\infty} \frac{\Lambda(n)}{n^{s}}=-\frac{\zeta^{\prime}}{\zeta}(s) \sim \frac{1}{s-1} \quad(\text { as } s \downarrow 1),
$$

and so it must be that $K= \pm 1$. From this we easily deduce that

$$
\lim _{s \downarrow 1}(s-1)\left(-\frac{L^{\prime}}{L}(s, \chi)\right)= \begin{cases}0 & \text { if } \chi \neq \chi_{0} \text { and } L(1, \chi) \neq 0, \\ -1 & \text { if } \chi \neq \chi_{0} \text { and } L(1, \chi)=0, \\ 1 & \text { if } \chi=\chi_{0} .\end{cases}
$$

The numbers on the right-hand side correspond precisely to the coefficients of $\log x$ in the estimates of Lemmas 4.5 and 4.6. This is not a coincidence! Indeed, that something like this should be true is a frequently useful principle in analytic number theory, which finds concrete embodiment in various socalled "Tauberian" theorems. See, e.g., [Ten95, §7.3].

## 6. Nonvanishing of $L(1, \chi)$ for real $\chi$

Lemma 4.8. Let $\chi$ be a real Dirichlet character modulo $m$. For every natural number $n$,

$$
\sum_{d \mid n} \chi(d) \geq \begin{cases}1 & \text { if } n \text { is a perfect square }, \\ 0 & \text { in any case } .\end{cases}
$$

Proof. Let $F(n):=\sum_{d \mid n} \chi(d)$. Since $\chi$ is multiplicative, $F$ is also multiplicative. Hence $F(n)=\prod_{p^{e} \| n} F\left(p^{e}\right)$. Since $\chi$ is real, we have $\chi(p)=0,1$, or -1 , so that

$$
F\left(p^{e}\right)=1+\chi(p)+\cdots+\chi\left(p^{e}\right)= \begin{cases}1 & \text { if } \chi(p)=0, \\ e+1 & \text { if } \chi(p)=1, \\ 0 & \text { if } \chi(p)=-1 \text { and } 2 \nmid e, \\ 1 & \text { if } \chi(p)=-1 \text { and } 2 \mid e .\end{cases}
$$

Since $F\left(p^{e}\right)$ is always nonnegative and $F\left(p^{e}\right) \geq 1$ when $e$ is even, the lemma follows.

Suppose now that $\chi$ is nontrivial. By partial summation,

$$
\begin{equation*}
\sum_{n \leq x} \frac{\chi(n)}{n}=\frac{S(x)}{x}+\int_{1}^{x} S(t) \frac{d t}{t^{2}}, \quad \text { where } \quad S(t):=\sum_{n \leq t} \chi(n) . \tag{4.26}
\end{equation*}
$$

Moreover, $S(t)$ is $O(1)$ (in fact, bounded by $\varphi(m)$ ). Multiplying (4.26) through by $x$ and recalling (Lemma 4.4) that

$$
L(1, \chi)-\sum_{n \leq x} \frac{\chi(n)}{n}=O\left(\frac{1}{x}\right)
$$

we find that for $x \geq 2$,

$$
\begin{aligned}
x L(1, \chi) & =\int_{1}^{x}\left(\sum_{n \leq t} \chi(n)\right) \frac{x}{t^{2}} d t+O(1) \\
& =\int_{1}^{x}\left(\sum_{n \leq t} \chi(n)\right)\left\lfloor\frac{x}{t}\right\rfloor \frac{1}{t} d t+O(\log x) \\
& =\int_{1}^{x}\left(\sum_{n \leq t} \chi(n) \sum_{a \leq x / t} 1\right) \frac{1}{t} d t+O(\log x) .
\end{aligned}
$$

This integral may be rewritten as

$$
\sum_{a n \leq x} \chi(n) \int_{n}^{x / a} \frac{1}{t} d t=\sum_{a n \leq x} \chi(n) \log \frac{x}{a n}=\sum_{N \leq x}\left(\sum_{d \mid N} \chi(d)\right) \log \frac{x}{N},
$$

which by Lemma 4.8 is bounded below by

$$
\sum_{M \leq \sqrt{x}} \log \frac{x}{M^{2}}=2 \sum_{M \leq \sqrt{x}} \log \frac{\sqrt{x}}{M} \geq 2 \log 2\left\lfloor\frac{\sqrt{x}}{2}\right\rfloor
$$

where the final bound comes from just considering those values of $M \leq$ $\sqrt{x} / 2$. Hence

$$
x L(1, \chi) \geq 2 \log 2\left\lfloor\frac{\sqrt{x}}{2}\right\rfloor+O(\log x) .
$$

The right-hand side of this inequality is positive for large enough $x$, which is only possible if $L(1, \chi)>0$.

## 7. Finishing up

Let $m$ be a positive integer and let $a$ be any integer coprime to $m$. We now know that $L(1, \chi)$ is nonvanishing for every nontrivial Dirichlet character $\chi$ modulo $m$. It follows from Lemma 4.5 that for every such $\chi$,

$$
\begin{equation*}
\sum_{n \leq x} \frac{\chi(n) \Lambda(n)}{n}=O(1) \tag{4.27}
\end{equation*}
$$

We record here also the result of Lemma 4.6 that

$$
\begin{equation*}
\sum_{n \leq x} \frac{\chi_{0}(n) \Lambda(n)}{n}=\log x+O(1) \tag{4.28}
\end{equation*}
$$

From the orthogonality relation (4.17), we see that

$$
\begin{aligned}
\sum_{\substack{n \leq x \\
n \equiv a(\bmod m)}} \frac{\Lambda(n)}{n} & =\frac{1}{\varphi(m)} \sum_{\chi} \bar{\chi}(a) \sum_{n \leq x} \frac{\chi(n) \Lambda(n)}{n} \\
& =\frac{1}{\varphi(m)} \bar{\chi}(a) \log x+O(1)=\frac{1}{\varphi(m)} \log x+O(1),
\end{aligned}
$$

since $\chi_{0}(a)=1$ (because $\left.\operatorname{gcd}(a, m)=1\right)$.
But already in the introduction we showed that

$$
\begin{equation*}
\sum_{\substack{p \leq x \\ p \equiv a(\bmod m)}} \frac{\log p}{p}=\sum_{\substack{n \leq x \\ n \equiv a(\bmod m)}} \frac{\Lambda(n)}{n}+O(1), \tag{4.30}
\end{equation*}
$$

with an absolute implied constant (see (4.6)). So from (4.29),

$$
\sum_{\substack{p \leq x \\ p \equiv a(\bmod m)}} \frac{\log p}{p}=\frac{1}{\varphi(m)} \log x+O(1) .
$$

This completes the proof of (4.1) in the general case.

## 8. Sums of three squares

Our goal in this section is to prove the following theorem of Legendre (see [Leg00, Troisième Partie]):

Theorem 4.9. A natural number $n$ can be written as the sum of three squares of integers if and only if $n$ does not have the form $4^{k}(8 l+7)$ for nonnegative integers $k$ and $l$.

We first dispense with the necessity half of Theorem 4.9.
Lemma 4.10. Suppose the positive integer $n$ is a sum of three squares of integers. Then $n$ is not of the form $4^{k}(8 l+7)$.

Proof. Suppose $n$ has the form $4^{k}(8 l+7)$ but that $n$ is a sum of three squares, say $n=x^{2}+y^{2}+z^{2}$. Since every square is either 0 or 1 modulo 4, if 4 divides $n$, we must have $x^{2} \equiv y^{2} \equiv z^{2} \equiv 0(\bmod 4)$, so that all of $x, y, z$ are even. Thus $n / 4=(x / 2)^{2}+(y / 2)^{2}+(z / 2)^{2}$ is also a sum of three squares. Applying this reasoning $k$ times, we eventually find that $8 l+7$ is a sum of three squares. But this is impossible, since the congruence $x^{2}+y^{2}+z^{2} \equiv 7(\bmod 8)$ has no solutions.

We can therefore focus our attention on the sufficiency portion of Theorem 4.9. Our proof of this requires another of Legendre's results (see [Leg00, Première Partie, §IV]), of independent interest:

Theorem 4.11 (Legendre). Suppose $a, b$, and $c$ are squarefree, pairwise coprime nonzero integers, not all of the same sign. In order that there exist a nonzero solution $(x, y, z) \in \mathbf{Z}^{3}$ to the equation

$$
\begin{equation*}
a x^{2}+b y^{2}+c z^{2}=0 \tag{4.31}
\end{equation*}
$$

it is necessary and sufficient that $-a b$ be a square modulo $c$, $-a c$ a square modulo $b$, and $-b c$ a square modulo $a$.

Before proving Theorem 4.11 we need the following simple but useful lemma, the proof of which is similar to an argument that appeared already in the proof of Lemma 2.3.

Say that two vectors with integer entries are congruent modulo $m$ if every entry in their difference is a multiple of $m$.

Lemma 4.12 (Brauer \& Reynolds [BR51]). Let $A=\left(a_{i j}\right)_{1 \leq i \leq r, 1 \leq j \leq s}$ be an $r \times s$ matrix with integer entries, and let $m$ be a natural number. Suppose that $\lambda_{1}, \ldots, \lambda_{s}$ are positive real numbers with $\lambda_{1} \cdots \lambda_{s}>m^{r}$. Then there is a nonzero column vector $\mathbf{x}=\left(x_{1}, \ldots, x_{s}\right)^{T}$ with integer entries satisfying $A \mathbf{x} \equiv \mathbf{0}(\bmod m)$ and having each $\left|x_{i}\right|<\lambda_{i}$.

Proof. For a real number $\lambda$, let $\lfloor\lambda \rrbracket$ be the greatest integer strictly less than $\lambda$. Let $\mathbf{x}=\left(x_{1}, \ldots, x_{s}\right)^{T}$ range over the $s \times 1$ vectors with integer entries $x_{i}$ satisfying $0 \leq x_{i}<\lambda_{i}$ for each $1 \leq i \leq s$. The number of such vectors $\mathbf{x}$ is $\prod_{i=1}^{s}\left(1+\left\lfloor\lambda_{i} \Downarrow\right) \geq \prod_{i=1}^{s} \lambda_{i}>m^{r}\right.$. This implies that there must be two distinct vectors of this type, say $\mathbf{x}_{1}$ and $\mathbf{x}_{2}$, for which $A \mathbf{x}_{1} \equiv A \mathbf{x}_{2}(\bmod m)$. The theorem follows with $\mathbf{x}:=\mathbf{x}_{1}-\mathbf{x}_{2}$.

Proof of Theorem 4.11. First we prove necessity. Suppose $(x, y, z)$ is a nonzero solution to (4.31). Dividing (4.31) by $\operatorname{gcd}(x, y, z)^{2}$, we can assume from the start that $\operatorname{gcd}(x, y, z)=1$. Considering (4.31) modulo $c$, we find that $a x^{2} \equiv-b y^{2}(\bmod c)$, so that

$$
\begin{equation*}
(a x)^{2} \equiv(-a b) y^{2}(\bmod c) \tag{4.32}
\end{equation*}
$$

Moreover, $y$ is invertible modulo $c$ : Otherwise, there is a prime $p$ dividing both $c$ and $y$. From (4.32), this $p$ divides $a x$; since $\operatorname{gcd}(a, c)=1$, it follows that $p$ divides $x$. But then $p^{2} \mid a x^{2}+b y^{2}=-c z^{2}$, and since $c$ is squarefree, we obtain that $p$ divides $z$. Thus $p$ divides $\operatorname{gcd}(x, y, z)$, a contradiction. So $y$ is invertible modulo $c$ and from (4.32) we get $\left(a x y^{-1}\right)^{2} \equiv-a b(\bmod c)$, so that $-a b$ is a square modulo $c$. The other necessary conditions are established similarly.

Now we turn to sufficiency. We claim that modulo $a b c$, the diagonal form $a x^{2}+b y^{2}+c z^{2}$ splits into linear factors. That is, there are integers $A_{1}, B_{1}, C_{1}$ and $A_{2}, B_{2}, C_{2}$ for which

$$
\begin{equation*}
a x^{2}+b y^{2}+c z^{2} \equiv\left(A_{1} x+B_{1} y+C_{1} z\right)\left(A_{2} x+B_{2} y+C_{2} z\right)(\bmod a b c) . \tag{4.33}
\end{equation*}
$$

By the Chinese remainder theorem, to prove the claim it is enough to show that a factorization of this type exists modulo each of $a, b$, and $c$. Suppose we first look modulo $a$. By hypothesis, we can choose an integer $u$ with $u^{2} \equiv$ $-b c(\bmod a)$. Then, using $b^{-1}$ to denote an integer with $b^{-1} b \equiv 1(\bmod a)$,

$$
\begin{aligned}
a x^{2}+b y^{2}+c z^{2} & \equiv b y^{2}+c z^{2} \equiv b^{-1}\left(b^{2} y^{2}+b c z^{2}\right) \\
& \equiv b^{-1}(b y-u z)(b y+u z) \equiv\left(y-b^{-1} u z\right)(b y+u z)(\bmod a)
\end{aligned}
$$

which is a factorization of the desired form. In exactly the same way we obtain factorizations modulo $b$ and modulo $c$, proving the claim.

Since $a, b$, and $c$ are not all of the same sign, we can assume $a, b>0$ and $c<0$. We can also assume $|a b c|>1$, since otherwise the theorem is trivial. Put $\lambda_{1}:=\sqrt{|b c|}, \lambda_{2}=\sqrt{|a c|}$, and $\lambda_{3}:=\sqrt{|a b|}$. Since either $|b c|,|a c|$, or $|a b|$ is squarefree and $>1$, not every $\lambda_{i}$ can be an integer. Pick one that is not, and increase it slightly, without changing $\left\lfloor\lambda_{i} \rrbracket\right.$. Then $\lambda_{1} \lambda_{2} \lambda_{3}>|a b c|$, and so from Lemma 4.12 (with $r=1$ and $s=3$ ), there are integers $x, y, z$, not all zero, with

$$
A_{1} x+B_{1} y+C_{1} z \equiv 0(\bmod a b c), \quad|x|<\sqrt{|b c|},|y|<\sqrt{|a c|},|z|<\sqrt{|a b|} .
$$

From (4.33), it follows that $a x^{2}+b y^{2}+c z^{2}$ is a multiple of $a b c$; moreover,

$$
-|a b c|<c z^{2} \leq a x^{2}+b y^{2}+c z^{2} \leq a x^{2}+b y^{2}<a|b c|+b|a c|=2|a b c| .
$$

So either $a x^{2}+b y^{2}+c z^{2}=0$ or $a x^{2}+b y^{2}+c z^{2}=|a b c|=-a b c$. In the first case we are done. In the second case,

$$
a x^{2}+b y^{2}+c\left(z^{2}+a b\right)=0 .
$$

Multiplying through by $z^{2}+a b$, we find
$0=\left(a x^{2}+b y^{2}\right)\left(z^{2}+a b\right)+c\left(z^{2}+a b\right)^{2}=a(x z+b y)^{2}+b(y z-a x)^{2}+c\left(z^{2}+a b\right)^{2}$.
Moreover, this is nontrivial since $z^{2}+a b>0$. So once again we are done.
The next lemma reduces our task to showing that a number $n$ meeting the conditions of Theorem 4.9 can be written as a sum of three squares of rational numbers.

Lemma 4.13. Suppose that the positive integer $n$ is the sum of three squares of rational numbers. Then $n$ is the sum of three squares of integers.

Proof (Aubry). If $n$ is a sum of three rational squares, then there is a point $\mathbf{a}=\left(a_{1}, a_{2}, a_{3}\right)$ with rational coordinates on the sphere $x^{2}+y^{2}+z^{2}=n$. Let $d$ be the least common denominator of $a_{1}, a_{2}, a_{3}$, so that

$$
A_{1}:=d a_{1}, A_{2}:=d a_{2}, A_{3}:=d a_{3} \quad \text { are integers, and } \operatorname{gcd}\left(A_{1}, A_{2}, A_{3}, d\right)=1 .
$$

Suppose that the rational point a is chosen so that $d$ is as small as possible. We shall show that $d=1$, so that a has integer coordinates, making $n$ a sum of three integer squares.

Suppose $d>1$. Let $\mathbf{a}^{\prime}=\left(a_{1}^{\prime}, a_{2}^{\prime}, a_{3}^{\prime}\right)$ be a point of $\mathbf{Z}^{3}$ closest to $\mathbf{a}$, so that

$$
\begin{equation*}
\left|a_{i}-a_{i}^{\prime}\right| \leq \frac{1}{2} \text { for each } 1 \leq i \leq 3, \quad \text { whence } \quad\left\|\mathbf{a}-\mathbf{a}^{\prime}\right\| \leq \frac{\sqrt{3}}{2}<1 \tag{4.34}
\end{equation*}
$$

Observe that

$$
\left\|\mathbf{a}-\mathbf{a}^{\prime}\right\|^{2}=\frac{1}{d^{2}} \sum_{i=1}^{3}\left(A_{i}-d a_{i}^{\prime}\right)^{2},
$$

while

$$
\begin{equation*}
\sum_{i=1}^{3}\left(A_{i}-d a_{i}^{\prime}\right)^{2} \equiv A_{1}^{2}+A_{2}^{2}+A_{3}^{2}=d^{2} n \equiv 0(\bmod d) \tag{4.35}
\end{equation*}
$$

By (4.34) and (4.35),

$$
\begin{equation*}
\left\|\mathbf{a}-\mathbf{a}^{\prime}\right\|^{2}=\frac{d^{\prime}}{d} \tag{4.36}
\end{equation*}
$$

for some $1 \leq d^{\prime}<d$. We shall exhibit a rational point on our sphere with (not necessarily least) common denominator $d^{\prime}$, contradicting the minimality of $d$.

This point will be the second intersection point of the line through a and $\mathbf{a}^{\prime}$ with the sphere $x^{2}+y^{2}+z^{2}=n$. Put $\mathbf{A}:=\left(A_{1}, A_{2}, A_{3}\right)$. Since $\mathbf{a}-\mathbf{a}^{\prime}=\left(\mathbf{A}-d \mathbf{a}^{\prime}\right) / d$, the line through $\mathbf{a}$ and $\mathbf{a}^{\prime}$ can be parameterized by a real parameter $\lambda$ as

$$
\mathbf{a}^{\prime}+\lambda\left(\mathbf{A}-d \mathbf{a}^{\prime}\right) .
$$

Setting the squared norm of this vector equal to $n$ gives the equation

$$
\left\|\mathbf{a}^{\prime}\right\|^{2}-n+2 \lambda\left(\mathbf{a}^{\prime} \cdot \mathbf{A}-d\left\|\mathbf{a}^{\prime}\right\|^{2}\right)+\left\|\mathbf{A}-d \mathbf{a}^{\prime}\right\|^{2} \lambda^{2}=0
$$

This is a quadratic equation in $\lambda$. We know already that $\lambda=1 / d$ is a root; this corresponds to the point a on the sphere. Since the roots multiply to

$$
\frac{\left\|\mathbf{a}^{\prime}\right\|^{2}-n}{\left\|\mathbf{A}-d \mathbf{a}^{\prime}\right\|^{2}}
$$

the root corresponding to the other intersection point is (by (4.36))

$$
\lambda=d \frac{\left\|\mathbf{a}^{\prime}\right\|^{2}-n}{\left\|A-d \mathbf{a}^{\prime}\right\|^{2}}=d \frac{\left\|\mathbf{a}^{\prime}\right\|^{2}-n}{d^{\prime} d}=\frac{\left\|\mathbf{a}^{\prime}\right\|^{2}-n}{d^{\prime}} .
$$

Thus $\lambda$ can be written as a fraction with denominator $d^{\prime}<d$, which implies that the same is true for the coordinates of the corresponding intersection point $\mathbf{a}^{\prime}+\lambda\left(\mathbf{A}-d \mathbf{a}^{\prime}\right)$.

We now complete the proof of sufficiency.

Lemma 4.14. Every positive integer not of the form $4^{k}(8 l+7)$ is a sum of three squares.

Proof. It is enough to prove that every squarefree positive integer $m \not \equiv$ $7(\bmod 8)$ is a sum of three squares. Indeed, suppose this special case is proven, and let $n$ be a positive integer not of the form $4^{k}(8 l+7)$. We can write $n=2^{2 k} a^{2} m$, where $k \geq 0, a$ is odd and $m$ is squarefree. The hypothesis on $n$ implies that

$$
m \equiv a^{2} m \not \equiv 7(\bmod 8) .
$$

Thus $m$ is a sum of three squares. Since $n$ is a square multiple of $m$, it follows that $n$ is also a sum of three squares.

To prove this special case we will construct a squarefree positive integer $r$ relatively prime to $m$ with the properties that
(i) $r$ is a sum of two integer squares,
(ii) $m$ is a square modulo $r$ and $-r$ is a square modulo $m$.

For this $r$, Legendre's theorem implies that there is are integers $x, y$, and $z$, not all zero, with

$$
m x^{2}-y^{2}-r z^{2}=0
$$

If $x=0$, then $y^{2}+r z^{2}=0$. But then also $y=z=0$, which is a contradiction. So $x \neq 0$, and we can divide through by $x^{2}$ to find

$$
m=(y / x)^{2}+r(z / x)^{2} .
$$

We are supposing that $r=r_{1}^{2}+r_{2}^{2}$ for integers $r_{1}$ and $r_{2}$, and thus

$$
m=(y / x)^{2}+\left(r_{1} z / x\right)^{2}+\left(r_{2} z / x\right)^{2} .
$$

So $m$ is a sum of three rational squares. By Lemma 4.13, $m$ is also a sum of three integer squares.

It remains to construct a suitable value of $r$. Write $m=2^{e} m_{1}$ where $e=0$ or 1 and $m_{1}=p_{1} \cdots p_{k}$ is odd. Put

$$
\beta:= \begin{cases}0 & \text { if } e=1, \text { or if } e=0 \text { and } m_{1} \equiv 1(\bmod 4), \\ 1 & \text { if } e=0 \text { and } m_{1} \equiv 3(\bmod 8) .\end{cases}
$$

Use Dirichlet's theorem to pick a prime $q$ with

$$
\begin{aligned}
&\left(\frac{q}{p_{i}}\right)=\left(\frac{-2^{\beta}}{p_{i}}\right) \text { for all } 1 \leq i \leq k, \\
& \text { and } q \equiv \begin{cases}1(\bmod 8) & \text { if } m_{1} \equiv 1(\bmod 4), \\
5(\bmod 8) & \text { if } m_{1} \equiv 3(\bmod 4) .\end{cases}
\end{aligned}
$$

(These conditions can be enforced by picking $q$ from a suitable residue class modulo $8 \prod p_{i}=8 m_{1}$.) We put $r:=2^{\beta} q$. Then classical results of Euler
show that $r$ can be written as a sum of two squares. Now $(q, m)=1$; moreover, $\beta>0$ only when $m$ is odd. Thus $r$ is coprime to $m$. Moreover, since $q \equiv 1(\bmod 4)$,

$$
\begin{aligned}
\left(\frac{m}{q}\right)=\left(\frac{2^{e}}{q}\right)\left(\frac{m_{1}}{q}\right) & =\left(\frac{2^{e}}{q}\right)\left(\frac{q}{m_{1}}\right)=\left(\frac{2^{e}}{q}\right) \prod_{i=1}^{k}\left(\frac{q}{p_{i}}\right) \\
& =\left(\frac{2^{e}}{q}\right) \prod_{i=1}^{k}\left(\frac{-2^{\beta}}{p_{i}}\right)=\left(\frac{2^{e}}{q}\right)\left(\frac{-2^{\beta}}{m_{1}}\right)=1 .
\end{aligned}
$$

(The last equality requires some checking of cases, depending on whether $e=0$ or 1 and whether $m_{1} \equiv 1$ or $3(\bmod 4)$.) Hence $m$ is a square modulo $q$. Since $m$ is trivially a square modulo $2^{\beta}$, we have by the Chinese remainder theorem that $m$ is a square modulo $r=2^{\beta} q$. Moreover,

$$
\left(\frac{-r}{p_{i}}\right)=\left(\frac{-2^{\beta} q}{p_{i}}\right)=\left(\frac{-2^{\beta}}{p_{i}}\right)\left(\frac{q}{p_{i}}\right)=\left(\frac{-2^{\beta}}{p_{i}}\right)^{2}=1 .
$$

By the Chinese remainder theorem, it follows that $-r$ is a square modulo $\prod p_{i}=m_{1}$. Since $-r$ is trivially a square modulo $2^{e}$, we have that $-r$ is also a square modulo $2^{e} m_{1}=m$, as desired.

## Notes

The proof of Dirichlet's theorem given here is a variant due to Gelfond [Gel56] of an argument of Shapiro ([Sha50], see also [Sha83, Chapter 9]). For the most part, our treatment follows that of Gelfond \& Linnik [GL66, $\S 3.2$, but the slick proof of the nonvanishing of $L(1, \chi)$ for real $\chi$ is due to Yanagisawa [Yan98]. A very different elementary proof of Dirichlet's theorem was given by Selberg [Sel49a]. An excellent presentation of the usual complex-analytic proof can be found in the textbook of Ireland \& Rosen [IR90, Chapter 16]. For a discussion of Dirichlet's original argument, and in particular his remarkable class number formula, see the beautiful text of Scharlau \& Opolka [SO85, Chapter 8].

For certain small moduli it is possible to prove Dirichlet's theorem by arguments analogous to those offered for Chebyshev's theorems in Chapter 3. See, e.g., Bang [Ban91, Ban37], Ricci [Ric33, Ric34], and Erdős [Erd35d]. Erdős's method, which is the most comprehensive, applies to any modulus $m$ for which $\sum_{p \nmid m, p<m} p^{-1}<1$. (This inequality has only finitely many solutions, the largest being $m=840$, as shown by Moree [Mor93].)

Shiu [Shi00] has established the following handsome strengthening of Dirichlet's theorem: If $a$ and $m$ are integers with $m>0$ and $\operatorname{gcd}(a, m)=1$, then for every $k \in \mathbf{N}$ the sequence of primes contains $k$ consecutive terms
each congruent to a modulo $m$. So, for example, there are $10^{100}$ consecutive primes each of which terminates in the decimal digit " 1 ".

Our proof of Theorem 4.9 characterizing sums of three squares is due to Wójcik [Wój72]. The proof of Legendre's Theorem 4.11 is based on the treatment of LeVeque [LeV96, Chapter 8]. From a modern perspective, Legendre's theorem is the first nontrivial case of the following important result of Hasse and Minkowski: If $Q$ is any quadratic form with rational coefficients, then $Q$ has a nontrivial zero over $\mathbf{Q}$ precisely when $Q$ has a nontrivial zero over $\mathbf{R}$ and every p-adic field $\mathbf{Q}_{p}$. The Hasse-Minkowski theorem can be used to give a quick proof of Theorem 4.9; see the appendix to [Ser73, Chapter IV].

In his Disquisitiones, Gauss determined the precise number of representations of an arbitrary natural number as a sum of three squares. For a natural number $n$, let $r_{3}(n)$ be the number of triples $(x, y, z) \in \mathbf{Z}^{3}$ with $x^{2}+$ $y^{2}+z^{2}=n$, and let $R_{3}(n)$ be the number of such triples with $\operatorname{gcd}(x, y, z)=1$. It is easy to see that $r_{3}(n)=\sum_{d^{2} \mid n} R_{3}\left(n / d^{2}\right)$. Gauss showed that $R_{3}(1)=6$, $R_{3}(2)=12, R_{3}(3)=8$, and for $n>3$,

$$
R_{3}(n)= \begin{cases}12 h(-4 n) & \text { if } n \equiv 1 \text { or } 2(\bmod 4),  \tag{4.37}\\ 24 h(-n) & \text { if } n \equiv 3(\bmod 8), \\ 0 & \text { otherwise }\end{cases}
$$

Here $h(D)$ is the number of classes of primitive binary quadratic forms of discriminant $D$ - explicitly, $h(D)$ is the number of solutions in integers $a$, $b$, and $c$ to $b^{2}-4 a c=D$, subject to the constraints that
$a>0$ and $c>0, \quad \operatorname{gcd}(a, b, c)=1, \quad$ and

$$
|b| \leq a \leq c, \text { with } b \geq 0 \text { if either }|b|=a \text { or } a=c
$$

(cf. [Gau86, Art. 291], [Ven70, Chapter 4, §16]). For $r_{3}(n)$ itself one has the following complicated explicit description: Let $T(n)$ denote the number of triples of positive integers $d, \delta, \delta^{\prime}$ where $d, \delta$, and $\delta^{\prime}$ are all odd, $d+\delta \equiv$ $0(\bmod 4)$, and $4 n+1=d \delta+(d+\delta \pm 2) \delta^{\prime}$ for some choice of sign. Then

$$
r_{3}(n)= \begin{cases}3 T(n) & \text { if } n \equiv 1 \operatorname{or} 2(\bmod 4)  \tag{4.38}\\ 2 T(n) & \text { if } n \equiv 3(\bmod 8) \\ r_{3}(n / 4) & \text { if } n \equiv 0(\bmod 4) \\ 0 & \text { if } n \equiv 7(\bmod 8)\end{cases}
$$

In Chapter XIII of the classic text of Uspensky \& Heaslet [UH39], one can find a a completely elementary proof of (4.38) based on certain remarkable identities of Liouville.

## Exercises

1. (Sylvester) Show that for complex $z$ with $|z|<1$,

$$
\sum_{\substack{n \geq 1 \\ p \mid n \Rightarrow p \equiv 3(\bmod 4)}} \mu(n) \frac{z^{n}}{1-z^{2 n}}=\sum_{\substack{m \geq 1 \\ p \mid m \Rightarrow p \equiv 1(\bmod 4)}} z^{m} .
$$

Suppose that there are only finitely many primes $p \equiv 3(\bmod 4)$. Setting $z=i y$ and letting $y$ tend to 1 from below, show that the left-hand side of this identity tends to a limit while the right-hand side "blows up" (has absolute value tending to infinity).
2. Let $\mathcal{P}$ be a set of primes. Suppose that

$$
\sum_{p \leq x, p \in \mathcal{P}} \frac{\log p}{p}=\kappa \log x+O_{\mathcal{P}}(1)
$$

for some constant $\kappa>0$ and every $x \geq 2$.
(a) Show that for some constant $D>1$, there are $\gg x / \log x$ elements of $\mathcal{P}$ in the interval $(x, D x]$ for every $x \geq 2$.
(b) Put $\pi_{\mathcal{P}}(x):=\#\{p \leq x: p \in \mathcal{P}\}$. Using the result of (a), show that $\pi_{\mathcal{P}}(x)>_{\mathcal{P}} x / \log x$ as $x \rightarrow \infty$.
(c) Show that if $\lim _{x \rightarrow \infty} \frac{\pi_{\mathcal{P}}(x)}{x / \log x}$ exists, then it equals $\kappa$.
3. Show that under the hypotheses of Exercise 2, there is a positive constant $C=C(\mathcal{P})$ for which

$$
\prod_{p \leq x, p \in \mathcal{P}}\left(1-\frac{1}{p}\right)=\frac{C}{(\log x)^{\kappa}}\left(1+O\left(\frac{1}{\log x}\right)\right)
$$

for $x \geq 2$. Here the implied constant may depend on $\mathcal{P}$.
Remark. When $\mathcal{P}$ is the set of primes $p \equiv a(\bmod m)$ (so that $\kappa=$ $1 / \varphi(m)$ ), Languasco \& Zaccagnini $[\mathbf{L Z 0 7}]$ have shown that $C$ is the positive solution to

$$
C^{\varphi(m)}=e^{-\gamma} \prod_{p}(1-1 / p)^{\alpha(p ; m, a)}
$$

where $\alpha(p ; m, a):=\varphi(m)-1$ if $p \equiv a(\bmod m)$ and $\alpha(p ; m, a)=-1$ otherwise.
4. Suppose that $\chi: \mathbf{Z} \rightarrow \mathbf{C}$ has the following three properties:
(i) $\chi$ is periodic with period $m$,
(ii) $\chi$ is completely multiplicative,
(iii) $\chi(n)=0$ if and only if $\operatorname{gcd}(n, m)>1$.

Show that $\chi$ is a Dirichlet character modulo $m$.
5. Let $G$ be a finite abelian group and let $\mathbf{C}[G]$ denote the space of functions $f: G \rightarrow \mathbf{C}$. For $\phi, \psi \in \mathbf{C}[G]$, define

$$
(\phi, \psi)=\frac{1}{|G|} \sum_{g \in G} \phi(g) \overline{\psi(g)}
$$

Show that this is a scalar product on $\mathbf{C}[G]$. Using (4.14) show that the characters of $G$ form an orthonormal basis for $\mathbf{C}[G]$. This explains the name "orthogonality relation".
6. Let $G$ be a finite abelian group of order $n$ with elements $g_{1}, g_{2}, \ldots, g_{n}$ and characters $\chi_{1}, \chi_{2}, \ldots, \chi_{n}$. Define the matrix

$$
M:=\left(\begin{array}{cccc}
\chi_{1}\left(g_{1}\right) & \chi_{1}\left(g_{2}\right) & \ldots & \chi_{1}\left(g_{n}\right) \\
\chi_{2}\left(g_{1}\right) & \chi_{2}\left(g_{2}\right) & \ldots & \chi_{2}\left(g_{n}\right) \\
\vdots & \vdots & \ddots & \vdots \\
\chi_{n}\left(g_{1}\right) & \chi_{n}\left(g_{2}\right) & \ldots & \chi_{n}\left(g_{n}\right)
\end{array}\right)
$$

Let $M^{*}$ denote the conjugate-transpose of $M$. Using (4.14), show that $M M^{*}=n I$, where $I$ is the $n \times n$ identity matrix. Linear algebra implies that $M^{*} M=n I$ as well. Deduce from this that (4.15) holds. That is, the first orthogonality relation implies the second.
7. (Sylvester, $[\mathbf{S y l 8 8}]$ ) Let $f$ be a nonnegative, multiplicative arithmetic function. Let $\chi$ be a nontrivial character modulo $m$, and define the arithmetic function $g$ by setting $g(n):=\sum_{d \mid n} \chi(d) f(n / d)$. Using only the convergence of $L(1, \chi)$ (and not its nonvanishing), prove that

$$
\left|\sum_{n \leq x} \frac{g(n)}{n}\right| \ll \prod_{p \leq x}\left(1+\frac{f(p)}{p}+\frac{f\left(p^{2}\right)}{p^{2}}+\cdots\right)
$$

where the implied constant depends at most on $m$.
8. (Continuation) Let $\chi$ be a nontrivial real Dirichlet character modulo $m$. Show that there is a unique choice of $f$ in the preceding exercise with the property that the induced function $g$ is identically 1 . Show, moreover, that this $f$ is nonnegative and multiplicative, and that for each prime $p$ and each $k \geq 1$ we have

$$
f\left(p^{k}\right)= \begin{cases}1 & \text { if } \chi(p)=0 \\ 0 & \text { if } \chi(p)=1 \\ 2 & \text { if } \chi(p)=-1\end{cases}
$$

Deduce from the preceding exercise that

$$
\sum_{\substack{p \leq x \\ \chi(p)=-1}} \frac{1}{p} \geq \frac{1}{2} \log \log x+O(1)
$$

As a special case, we see that the sum of the reciprocals of the primes from the residue class $3 \bmod 4$ diverges at least as fast as $\frac{1}{2} \log \log x$.
9. (Mertens $[\operatorname{Mer} 97])$ Suppose that $a$ and $m$ are integers with $m>0$ and $\operatorname{gcd}(a, m)=1$.
(a) Show that if $\chi$ is a character modulo $m$ and $x \geq 4$, then

$$
\sum_{d \leq x} \frac{\chi(d) \Lambda(d)}{d}= \begin{cases}\log x+O\left(1+\sum_{p \mid m} \frac{\log p}{p-1}\right) & \text { if } \chi=\chi_{0}, \\ O\left(m|L(1, \chi)|^{-1}\right) & \text { if } \chi \neq \chi_{0}\end{cases}
$$

where the implied constants are absolute.
(b) Put $M=\sum_{\chi \neq \chi_{0}}|L(1, \chi)|^{-1}$, where the sum is over all nontrivial Dirichlet characters modulo $m$. Show that

$$
\sum_{p \leq x} \frac{\log p}{p}=\frac{1}{\varphi(q)} \log x+O(M+1),
$$

again with an absolute implied constant.
(c) By splitting the sum defining $L(1, \chi)$ at $n=m$, show that $L(1, \chi) \ll$ $\log m$ for each nontrivial character $\chi$, so that $M \gg \varphi(m) / \log m$. (Again, both implied constants are supposed to be absolute.)
(d) Deduce that there is an absolute positive constant $C$ with the property that for every $x \geq 4$, there is a prime $p \equiv a(\bmod m)$ in the interval $[x, x \exp (C m M)]$.
Remark. Let $p(m, a)$ be the least prime $p \equiv a(\bmod m)$. From (d) we have that $p(m, a) \ll \exp (C m M)$. Unfortunately, from (c) this upper bound is quite large, at least $\exp \left(m^{2+o(1)}\right)$. See Révész [Rév80] for an elementary proof that $p(m, a) \ll \exp \left(c m(\log m)^{11}\right)$ for an absolute constant $c>0$.

A deep result of Linnik asserts that $p(m, a) \ll m^{L}$ for an absolute constant $L$. Heath-Brown has shown [HB92] that one may take $L=5.5$. The so-called Extended Riemann Hypothesis (which asserts that all the "nontrivial" zeros of the functions $L(s, \chi)$ lie on the line $\Re(s)=1 / 2$ ) would imply that $p(m, a)<2 m^{2}(\log m)^{2}$ (see [BS96]).
10. (Sierpiński [Sie62]) Suppose that $a$ and $m$ are coprime integers with $m>0$. Prove that for every $s \in \mathbf{N}$, there are infinitely many natural numbers $n \equiv a(\bmod m)$ with exactly $s$ prime divisors (counted with multiplicity).
11. (Schinzel [Sch59]) Prove that there are no congruence obstructions to the Goldbach conjecture. That is, show that if $n$ is an even integer and $m$ is a (positive) modulus, then the congruence $n \equiv p+q(\bmod m)$ is always solvable in primes $p$ and $q$.
12. (Sierpiński [Sie48]) Prove that for each $M \in \mathbf{N}$, there are infinitely many primes $p$ for which all of $p \pm i, i=1,2, \ldots, M$, are composite.
13. (Powell, Israel [ $\mathbf{I s r} 83]$ ) Let $m$ and $n$ be natural numbers with $m>1$. Show that if $(m, n) \neq(2,1)$, then $m^{p}-n$ is composite for infinitely many primes $p$.
14. (Newman [New97]; see also Aldaz et al. [ABGU01]) Dov Jarden, in his book Recurring Sequences (1973), observed that $\varphi(30 n+1)>\varphi(30 n)$ for all $n \leq 10,000$.

Prove that contrary to what one might expect from the computational evidence, the reverse inequality,

$$
\begin{equation*}
\varphi(30 n+1)<\varphi(30 n), \tag{4.39}
\end{equation*}
$$

holds for infinitely many $n$. Hint: Consider large primes $n$ for which $30 n+1$ has many small prime factors.

Remark. The smallest solution to (4.39), which has over 1000 decimal digits, has been given explicitly by Martin [Mar99].
15. This exercise illustrates the utility of (4.1) as an equidistribution statement. Define $n^{\diamond}$ as that portion of $n!$ composed of primes congruent to $3(\bmod 4)$, i.e., $n^{\diamond}:=\prod_{p^{k} \| n!, p \equiv 3(\bmod 4)} p^{k}$.
(a) Using (4.1), show that $\log n^{\diamond} \sim \frac{1}{2} \log n$ !. Hint: First show that if $p$ is prime, then $p^{k} \| n$ ! for $k=\sum_{i \geq 1}\left\lfloor n / p^{i}\right\rfloor$.
(b) Suppose that $n$ and $y$ are positive integers with $n!+1=y^{8}$. Using the factorization

$$
n!=y^{8}-1=\left(\left(y^{4}+1\right)\left(y^{2}+1\right)\right)\left(y^{2}-1\right),
$$

prove that $n^{\diamond} \leq y^{2}-1 \leq(n!)^{1 / 4}$. Deduce from part (a) that the equation $n!+1=y^{8}$ has only finitely many solutions.
(c) Show that the equation $n!+1=x^{p}$ has at most finitely many solutions ( $n, x$ ) for each fixed odd prime $p$.
In combination with the result of (b), this shows that $n!+1=x^{m}$ has only finitely many solutions for each positive integer $m>1$ except possibly for $m=2$ and $m=4$.

Remark. It has been shown that $n!+1=y^{m}$ has no solutions for any $m>2$. See $[\mathbf{E O 3 7}]$ for the case $m \neq 4$ and $[\mathbf{P S 7 3}]$ for the case $m=4$. When $m=2$, an 1885 conjecture of Brocard asserts that the only solutions correspond to $n=4,5$ and 7 , but even showing there are at most finitely many solutions remains open.
16. (Continuation; Dąbrowski [Dạb96]) Show that if $A \in \mathbf{Z}$ is not a perfect square, then the equation $n!+A=y^{2}$ has only finitely many integral solutions.
17. (Chebyshev, Nagell [Nag22, §1]) For each $x \geq 1$, put

$$
N_{x}:=\prod_{n \leq x}\left(n^{2}+1\right)
$$

(a) Show that $\log N_{x}=2 x \log x+O(x)$.
(b) For each prime $p$, define $e(p, x)$ by the relation $p^{e(p, x)} \| N_{x}$. Show that

$$
e(p, x) \leq \begin{cases}x / 2+O(1) & \text { if } p=2 \\ 2 x /(p-1)+O(\log x / \log p) & \text { if } p \equiv 1(\bmod 4) \\ 0 & \text { if } p \equiv 3(\bmod 4)\end{cases}
$$

(c) Show that there is a constant $c>0$ with the property that the largest prime factor $p_{x}$ of $N_{x}$ satisfies $p_{x}>c x \log x$ for all large $x$. Conclude that there are infinitely many $n \in \mathbf{N}$ for which $n^{2}+1$ has a prime factor $>c n \log n$. This can be considered an approximation to the conjecture that $n^{2}+1$ is prime infinitely often.

Remark. Deshouillers \& Iwaniec [DI82], building on earlier work of Hooley, have shown that $p_{x}>x^{\theta+o(1)}$ infinitely often, where $\theta$ is a constant slightly larger than $6 / 5$.
18. (Continuation; Cilleruelo [Cil08]) Show that if $e(p, x) \geq 2$, then $p \leq 2 x$. Deduce that $N_{x}$ assumes only finitely many squarefull values for $x \geq 1$. (With a bit more work, it can be shown that $10^{2}=\left(1^{2}+1\right)\left(2^{2}+1\right)\left(3^{2}+1\right)$ is the only such value.)
19. Show that $2 x+5 y+7$ and $3 x-2 y+1$ are simultaneously prime for infinitely many pairs $(x, y) \in \mathbf{Z}^{2}$. For the general theorem of which this is a special case, see [Tul83].
20. Let $\Phi(N)$ denote the number of Farey fractions of order $N$; in other words, $\Phi(N)$ is the number of reduced fractions $0 \leq \frac{a}{b} \leq 1$ with denominator $b \leq N$. It is not hard to see that

$$
\Phi(N)=1+\sum_{k=1}^{N} \varphi(k) .
$$

The first few values of $\Phi$ are

$$
2,3,5,7,11,13,19,23,29,33,43,47,59,65,73,81,97, \ldots
$$

Probably there are infinitely many primes in the sequence $\{\Phi(N)\}_{N=1}^{\infty}$, but this is presumably very hard. In this exercise we outline a proof (due to C. Pomerance) that the sequence $\{\Phi(N)\}_{N=1}^{\infty}$ hits every residue class modulo 3 infinitely often. In particular, there are infinitely many composite terms in this sequence.
(a) Let $\chi$ be the nontrivial Dirichlet character modulo 3. For real values of $x$, put $D(x):=\sum_{n \leq x} \chi(\varphi(n))$. Show that if some residue class modulo 3 contains only finitely many of the terms $\Phi(N)$, then $D(x)$ is absolutely bounded.
(b) Put $L(s):=\sum_{n=1}^{\infty} \frac{\chi(\varphi(n))}{n^{s}}$. Show that for real $s>1$, one has

$$
L(s)=\left(1-\frac{1}{3^{s}}\right) \prod_{p \equiv 2(\bmod 3)}\left(1+\frac{1}{p^{s}+1}\right) .
$$

(c) Conclude from (b) and the divergence of the series $\sum_{p \equiv 2(\bmod 3)} \frac{1}{p}$ that $L(s)$ tends to infinity as $s$ tends to 1 from above.
(d) Use the result of (c) to show that for each $\delta<1$, there are arbitrarily large values of $x$ with $D(x)>x^{\delta}$. In particular, $D(x)$ is not absolutely bounded.

Remark. The author does not know any proof of the analogous result for residue classes modulo 5 , or even a proof that 5 divides infinitely many of the terms $\Phi(N)$.
21. Let $p$ be a prime, and let $\zeta_{p}:=e^{2 \pi i / p}$, so that $\zeta_{p}$ is a complex primitive $p$ th root of unity. For each nontrivial character $\chi$ modulo $p$, define the Gauss sum $\tau(\chi)$ by setting

$$
\tau(\chi):=\sum_{n=1}^{p-1} \chi(n) \zeta_{p}^{n}
$$

(cf. Exercise 2.10, where certain analogous quantities were defined in positive characteristic).
(a) Show that $\tau(\chi) \overline{\tau(\chi)}=p$, so that $|\tau(\chi)|=\sqrt{p}$. You may wish to consult the hint to Exercise 2.10(a).
(b) For each integer $a$, define $\tau_{a}(\chi):=\sum_{n=1}^{p-1} \chi(n) \zeta_{p}^{a n}$. (Thus $\tau(\chi)=$ $\tau_{1}(\chi)$.) Show that $\tau_{a}(\chi)=\bar{\chi}(a) \tau(\chi)$.
(c) Deduce from the result of (b) that for each natural number $N$ and each nontrivial character $\chi$,

$$
\tau(\bar{\chi}) \sum_{a \leq N} \chi(a)=\sum_{n=1}^{p-1} \bar{\chi}(n) \sum_{a \leq N} \zeta_{p}^{a n} .
$$

Show that the right-hand inner sum has absolute value $\frac{\left|\sin \frac{\pi N n}{n}\right|}{\left|\sin \frac{\pi n}{p}\right|}$.
(d) Check that if $|\theta| \leq 1 / 2$, then $|\sin \pi \theta| \geq 2|\theta|$. Use this to prove the Pólya-Vinogradov inequality: For every $N$,

$$
\sum_{a \leq N} \chi(a)<\sqrt{p} \log p .
$$

22. (Continuation) Let $p$ be a prime.
(a) Suppose that $d$ divides $p-1$. Show that for $a$ coprime to $p$,
$\frac{1}{d} \sum_{\chi^{d}=\chi_{0}} \chi(a)= \begin{cases}1 & \text { if } a \text { is a } d \text { th power residue modulo } p, \\ 0 & \text { otherwise } .\end{cases}$
Here the sum on the left is extended over all characters modulo $p$ whose $d$ th power is the trivial character. Hint: The characters with $\chi^{d}=\chi_{0}$ can be identified in a natural way with the characters on the group $\mathbf{F}_{p}^{\times} /\left(\mathbf{F}_{p}^{\times}\right)^{d}$.
(b) Deduce from the Pólya-Vinogradov inequality that if $I$ is a finite interval of measure $\mu(I)$, then the number of $d$ th power residues in $I$ is $\mu(I) / d+O\left(p^{1 / 2} \log p\right)$. (Thus if $\mu(I)$ is significantly larger than $d p^{1 / 2} \log p$, then $I$ contains roughly the expected number of $d$ th power residues.)
(c) Show that for $a$ coprime to $p$,

$$
\sum_{e \mid p-1} \frac{\mu(e)}{e} \sum_{\chi^{e}=\chi_{0}} \chi(a)= \begin{cases}1 & \text { if } a \text { is a primitive root modulo } p \\ 0 & \text { otherwise } .\end{cases}
$$

(d) Prove that for each finite interval $I$, the number of primitive roots contained in $I$ is

$$
\mu(I) \frac{\varphi(p-1)}{p-1}+O\left(2^{\omega(p-1)} p^{1 / 2} \log p\right) .
$$

As a special case, conclude that if we let $g(p)$ denote the least positive primitive root modulo $p$, then for each $\epsilon>0$, we have $g(p) \ll_{\epsilon} p^{1 / 2+\epsilon}$.
Remark. Burgess [Bur62] has shown that $g(p) \ll_{\epsilon} p^{1 / 4+\epsilon}$ for each $\epsilon>0$, which is the best known unconditional upper bound. The Generalized Riemann Hypothesis implies (see [Sho92]) that $g(p) \ll(\log p)^{6}$, and it is conjectured that $g(p)<_{\epsilon}(\log p)^{1+\epsilon}$ for each $\epsilon>0$. In the opposite direction, there are infinitely many primes $p$ for which $g(p) \gg$ $\log p \log \log p ;$ in fact, the same lower bound holds for the least positive quadratic nonresidue [GR90].
23. By imitating the proof of Lemma 4.13, show that if the positive integer $n$ is a sum of two squares of rational numbers, then it is a sum of two squares of integers. Use this and Theorem 4.11 to show that every prime $p \equiv 1(\bmod 4)$ is a sum of two squares.
24. (Gauss [Gra84, Entry 18], [Gau86, Art. 293]) Show that every nonnegative integer $n$ can be written as a sum of three triangular numbers. (Here a triangular number is a number of the form $k(k+1) / 2$, where $k$ is a nonnegative integer.)


Figure 1. Pictorial representation of the first few nonzero $(m+2)$-gonal numbers when $m=3$. In this case the $j$ th step in the construction corresponds to adding $1+3 j$ dots.
25. Prove that the set of positive integers expressible as a sum of three squares has asymptotic density $5 / 6$.
26. (Turski [Tur33]) Prove that every positive integer is the sum of at most 10 odd squares and that infinitely many require 10.
27. Show that every nonnegative integer $n$ can be written as the sum of four squares of integers, where one of the integers belongs to the set $\{0\} \cup\left\{2^{k}: k=0,1,2, \ldots\right\}$.
28. (Sierpiński; extracted from [Sie88, Chapter XI]) Let $m$ be an odd positive integer.
(a) Prove that $m$ can be written as a sum of four squares of integers with two of the integers equal.
(b) Prove that $m$ can be written as a sum of four squares of integers with two of the integers consecutive.
Suggestions: For (a), write $2 m=x^{2}+y^{2}+z^{2}$. Show that we can assume that $x$ and $y$ are odd while $z$ is even. Verify that $m=((x+y) / 2)^{2}+$ $((x-y) / 2)^{2}+(z / 2)^{2}+(z / 2)^{2}$.

For (b), start by writing $2 m-1=x^{2}+y^{2}+z^{2}$. Show that after a rearrangement we can assume $x$ and $y$ are even while $z=2 c+1$ is odd. Now use the identity $c^{2}+(c+1)^{2}=\frac{1}{2}\left((2 c+1)^{2}+1\right)$.
29. For each natural number $m$, the sequence of $(m+2)$-gonal numbers is the sequence with $k$ th term

$$
p_{m}(k):=\sum_{0 \leq j<k}(1+m j)=\frac{m k^{2}-(m-2) k}{2},
$$

indexed starting at $k=0$. Figure 1 explains the geometric origin of the terminology. When $m=1$ we recover the triangular numbers of Exercise 24 , and when $m=2$ we recover the familiar sequence of square numbers. Fermat recorded the following claim in his copy of Diophantus's Arithmetica:

Every number is either a triangular number or the sum of two or three triangular numbers; every number is a square or the sum of two, three, or four squares; every number is a pentagonal number or the sum of two, three, four or five pentagonal numbers; and so on ad infinitum, for hexagons, heptagons, and any polygons whatever ... The proof, which depends on many various and abstruse mysteries of numbers, I cannot give here...
This statement is true; however, no record survives of Fermat's proof. The first published proof of the polygonal number theorem, as it has come to be called, is due to Cauchy [Cau15]. The argument is technical for uninteresting reasons, and we do not give it here. We can, however, sketch a proof of the following related theorem of Legendre (see [Leg00, Sixième Partie, §II]):
$\star$ Theorem 4.15. Fix a natural number $m$. If $m$ is odd, then every large enough natural number $n$ is a sum of four ( $m+2$ )-gonal numbers. If $m$ is even, then every large enough $n$ is a sum of five polygonal numbers of order $m+2$, one of which is either 0 or 1 .

Our sketch is based on [Nat96, Chapter 1], which also contains a proof of the polygonal number theorem.

The first step towards Theorem 4.15 is proving "Cauchy's lemma": If $a$ and $b$ are odd positive integers with $3 a \leq b^{2} \leq 4 a$, then there are nonnegative integers $s, t, u$, and $v$ with

$$
s+t+u+v=b \quad \text { and } \quad s^{2}+t^{2}+u^{2}+v^{2}=a .
$$

Proceed as follows:
(a) Deduce from Theorem 4.9 that we can write $4 a-b^{2}=x^{2}+y^{2}+z^{2}$ where $x, y$, and $z$ are odd integers.
(b) Show that one can choose the signs of $x, y$, and $z$ in (a) so that

$$
\begin{aligned}
& s:=\frac{b+x+y+z}{4}, \quad t:=\frac{b+x-y-z}{4}, \\
& \\
& \quad u:=\frac{b-x+y-z}{4}, \quad v:=\frac{b-x-y+z}{4}
\end{aligned}
$$

are all integers. Check that $s+t+u+v=b$ and $s^{2}+t^{2}+u^{2}+v^{2}=a$.
(c) Show that $|x|+|y|+|z| \leq b$ and conclude that each of $s, t, u$, and $v$ is nonnegative.
30. (Continuation) We can suppose for the proof of Theorem 4.15 that $m>$ 1, by Exercise 24. Suppose also that $m$ is odd.
(a) Show that if $n \geq 120 m^{3}$, then there is an odd natural number $b$ with $b \equiv n(\bmod m)$ and $\sqrt{7 n / m} \leq b \leq \sqrt{8 n / m}$.
(b) With $b$ as in (a), put $a:=2 \frac{n-b}{m}+b$. (Thus $a \equiv b \equiv 1(\bmod 2)$.) Show that (still under the assumption $n \geq 120 m^{3}$ ) $3 a \leq b^{2} \leq 4 a$.
(c) Choose $s, t, u, v$ as in Cauchy's lemma to correspond to the integers $a$ and $b$. Show that

$$
n=p_{m}(s)+p_{m}(t)+p_{m}(u)+p_{m}(v) .
$$

Thus every natural number $n \geq 120 m^{3}$ is a sum of four polygonal numbers of order $m+2$.
31. (Continuation) Suppose now that $m$ is even. Show that (a)-(c) of the preceding exercise remain correct if we make the extra assumption that $n$ is odd. Now complete the proof of Theorem 4.15.

Remark. Legendre (cf. [Nat87b]) actually proved a little bit more than Theorem 4.15: He showed that if $4 \mid m$, then every large enough $n$ is a sum of four polygonal numbers of order $m+2$, while if $2 \| m$, then every large enough $n \equiv 2(\bmod 4)$ is such a sum. For hexagonal numbers (corresponding to $m=4$ ), Duke (see [Duk97]) has shown that actually three such numbers suffice for large $n$; this is easily seen to be best-possible in this case. Some recent results and conjectures in the spirit of Cauchy's polygonal number theorem are discussed in [Sun] (see also [Sun07, GPS07, OS09]).

## Interlude: A Proof of the Hilbert-Waring Theorem

Every integer is a cube or the sum of at most nine cubes; every integer is also the square of a square, or the sum of up to nineteen such, and so forth. - E. Waring [War91]

It would hardly be possible for me to exaggerate the admiration which I feel for the solution of this historic problem... Within the limits which it has set for itself, it is absolutely and triumphantly successful, and it stands with the work of Hadamard and de la Vallée-Poussin, in the theory of primes, as one of the landmarks in the modern history of the theory of numbers. - G. H. Hardy [Har20] on Hilbert's solution of Waring's problem

## 1. Introduction

Fix an integer $k \geq 2$. Then every natural number $n$ can be written as a sum of nonnegative $k$ th powers, since trivially

$$
n=\overbrace{1^{k}+1^{k}+\cdots+1^{k}}^{n \text { terms }} .
$$

Of course this way of writing $n$ as a sum of $k$ th powers is usually vastly inefficient. Write $g(k ; n)$ for the minimal number of nonnegative $k$ th powers needed to additively represent $n$. (So, for example, $g(2 ; 5)=2$, since $5=$
$2^{2}+1^{2}$ and 5 is not a perfect square.) Let $g(k)$ be the supremum of the numbers $g(k ; n)$ as $n$ ranges over the set of natural numbers. In 1770, Waring asserted that $g(k)<\infty$ for every fixed $k$, and he conjectured that $g(3) \leq 9$ and $g(4) \leq 19$. (Presumably he believed equality to hold in both cases.)

Waring's claims have engaged the energies of mathematicians throughout the intervening centuries: The same year that Waring announced these conjectures, Lagrange proved his "four squares theorem" asserting that $g(2)=4$. In 1909, Wieferich [Wie09] proved that $g(3)=9$ (modulo a gap later filled by Kempner [Kem12]). Finally, in 1986, Balasubramanian et al. [BDD86a, BDD86b] succeeded in showing that $g(4)=19$. As described in the notes to this chapter, the precise value of $g(k)$ is now known for every $k$.

Our goals for this chapter are rather modest. We will not determine the exact value of $g(k)$ for even a single value of $k>2$. Instead, we describe what seems to be the simplest known proof of Waring's claim that all the numbers $g(k)$ are finite:

Theorem 5.1. $g(k)<\infty$ for each fixed $k$.
Theorem 5.1 was first established by Hilbert in 1909 [Hil09]. The proof presented here is a variant due to Dress [Dre71, Dre72a] of Hilbert's argument.

## 2. Proof of the Hilbert-Waring theorem (Theorem 5.1)

Fundamental to the proof of Theorem 5.1 is the following lemma which guarantees the existence of polynomial identities of a convenient shape:

Lemma 5.2 (Hilbert-Dress identities). Let $k \in \mathbf{N}$, and put $N:=\binom{2 k+4}{4}$. There is a formal identity in indeterminates $X_{1}, \ldots, X_{5}$ of the form

$$
\begin{equation*}
M\left(X_{1}^{2}+\cdots+X_{5}^{2}\right)^{k}=\sum_{i=0}^{N} m_{i}\left(a_{i 1} X_{1}+\cdots+a_{i 5} X_{5}\right)^{2 k}+m_{N+1} X_{5}^{2 k}, \tag{5.1}
\end{equation*}
$$

where $M$ and $m_{N+1}$ are positive integers, the $m_{i}(0 \leq i \leq N)$ are nonnegative integers, and the $a_{i j}(0 \leq i \leq N, 1 \leq j \leq 5)$ are integers.

The rather complicated proof of Lemma 5.2 is deferred to $\S 3$. Lemma 5.2 has the following important consequence:

Lemma 5.3. Fix a natural number $k$ and fix a corresponding identity of the form (5.1). Then with $M$ as in (5.1), one can find a natural number $Q$ with the following property: For every nonnegative integer $l$ and every integer $x$
with $|x| \leq \sqrt{l}$, we have

$$
M l^{k}=x^{2 k}+\sum_{h=1}^{Q} u_{h}^{2 k},
$$

for some integers $u_{1}, u_{2}, \ldots, u_{Q}$.
Thus, if we fix an $M$ as in (5.1), then each number of the form $M l^{k}$ can be written as the sum of $O_{k}(1)(2 k)$ th powers, one of which can be arbitrarily prescribed subject to a constraint on its size.

Proof. If $|x| \leq \sqrt{l}$, then by Lagrange's result on sums of four squares, we may write $l-x^{2}=x_{1}^{2}+x_{2}^{2}+x_{3}^{2}+x_{4}^{2}$ where $x_{1}, \ldots, x_{4} \in \mathbf{Z}$. Put $x_{5}:=x$. Evaluating both sides of (5.1) with $X_{i}=x_{i}$ for $1 \leq i \leq 5$, we find that

$$
\begin{aligned}
M l^{k}= & x^{2 k}+\overbrace{x^{2 k}+\cdots+x^{2 k}}^{m_{N+1}-1 \text { terms }} \\
& +\sum_{i=0}^{N}(\overbrace{\left(a_{i 1} x_{1}+\cdots+a_{i 5} x_{5}\right)^{2 k}+\cdots+\left(a_{i 1} x_{1}+\cdots+a_{i 5} x_{5}\right)^{2 k}}^{m_{i} \text { terms }}) .
\end{aligned}
$$

This proves the lemma with $Q:=m_{N+1}-1+\sum_{i=0}^{N} m_{i}$.
The next lemma guarantees the existence of another family of polynomial identities; these identities have long been well-known, but their use in the proof of Theorem 5.1 is due to Dress.

Lemma 5.4. For every natural number $k$, there is a formal identity in the indeterminate $T$ of the shape

$$
\sum_{i=1}^{R}\left(T+a_{i}\right)^{2 k}-\sum_{j=1}^{R}\left(T+a_{j}^{\prime}\right)^{2 k}=A T+B
$$

Here $R$ and $A$ are natural numbers and $B, a_{1}, \ldots, a_{R}, a_{1}^{\prime}, \ldots, a_{R}^{\prime}$ are integers. In fact, one can take $R=2^{2 k-2}$ and $A=(2 k)$ !.

Proof. The proof uses two easily-verified properties of the forward difference operator $\Delta: \mathbf{Z}[T] \rightarrow \mathbf{Z}[T]$, defined by

$$
(\Delta F)(T):=F(T+1)-F(T) .
$$

First, if $a_{n} T^{n}$ is the leading term of $F(T)$, where $n>0$, then $(\Delta F)(T)$ has leading term $n a_{n} T^{n-1}$. The second property concerns repeated application
of $\Delta$. Write $\Delta^{r}$ for $\Delta \circ \cdots \circ \Delta(r$ times $)$. Then for each natural number $r$ and each $F(T) \in \mathbf{Z}[T]$,

$$
\begin{equation*}
\left(\Delta^{r} F\right)(T)=\sum_{j=0}^{r}\binom{r}{j}(-1)^{r-j} F(T+j) \tag{5.2}
\end{equation*}
$$

Now take $F(T):=T^{2 k}$. Applying the first property $2 k-1$ times, we find that $\left(\Delta^{2 k-1} F\right)(T)=(2 k)!T+B$ for some integer $B$. So by the second property (with $r=2 k-1$ ), we conclude that

$$
(2 k)!T+B=\sum_{\substack{0 \leq j \leq 2 k-1 \\ 2 \nmid j}}\binom{2 k-1}{j}(T+j)^{2 k}-\sum_{\substack{0 \leq j \leq 2 k-1 \\ 2 \mid j}}\binom{2 k-1}{j}(T+j)^{2 k}
$$

Since

$$
\sum_{0 \leq j \leq 2 k-1}^{2 \nmid j}\binom{2 k-1}{j}=\sum_{0 \leq j \leq 2 k-1}^{2 \mid j} 0\binom{2 k-1}{j}=\frac{1}{2} \sum_{0 \leq j \leq 2 k-1}\binom{2 k-1}{j}=2^{2 k-2}
$$

the lemma follows with $R=2^{2 k-2}$ and $A=(2 k)$ !.
The last result we need is a simple lemma concerning how closely one can approximate a nonnegative real number by a sum of $k$ th powers of nonnegative integers:

Lemma 5.5. Let $k$ be a natural number and put $\kappa:=(k-1) / k$. Then for each $x \geq 0$, we have

$$
0 \leq x-\left\lfloor\left. x^{1 / k}\right|^{k} \leq k x^{\kappa}\right.
$$

Consequently, for all $x \geq 0$ and $t \in \mathbf{N}$, there are nonnegative integers $z_{1}, \ldots, z_{t}$ for which

$$
x-z_{1}^{k}-z_{2}^{k}-\cdots-z_{t}^{k} \ll_{k, t} x^{\kappa^{t}}
$$

Proof. By the mean value theorem, there is an $x^{\prime} \in\left(\left\lfloor x^{1 / k}\right\rfloor, x^{1 / k}\right)$ for which

$$
0 \leq x-\left\lfloor\left. x^{1 / k}\right|^{k}=\left.\frac{d}{d x}\left(x^{k}\right)\right|_{x=x^{\prime}}=k\left(x^{\prime}\right)^{k-1} \leq k x^{(k-1) / k}=k x^{\kappa}\right.
$$

Iterating this observation gives the lemma.
Proof of Theorem 5.1. Fix a natural number $k$. We wish to show that $g(k ; m)$ is bounded independently of $m$. Clearly it is enough to prove this for large $m$. To this end, fix an $R$ as in Lemma 5.4 and fix $M$ as in Lemmas 5.2 and 5.3. (Thus $R$ and $M$ depend entirely on $k$.)

Let $m$ be a large natural number which we seek to write as a sum of nonnegative $k$ th powers, and let $l^{k}$ be the largest $k$ th power not exceeding $m / R M$. If $m$ is sufficiently large, then

$$
\begin{equation*}
\frac{1}{2}\left(\frac{m}{R M}\right)^{1 / k} \leq l \leq\left(\frac{m}{R M}\right)^{1 / k} \tag{5.3}
\end{equation*}
$$

Moreover, by Lemma 5.5 with $x=m / R M$,

$$
m=R M l^{k}+r_{1}, \quad \text { where } \quad 0 \leq r_{1} \leq k R M\left(\frac{m}{R M}\right)^{(k-1) / k}
$$

With $\kappa:=(k-1) / k$, let $t$ be the least natural number for which $\kappa^{t}<\frac{1}{2 k}$. By Lemma 5.5 (with $x=r_{1}$ ) we can write

$$
\begin{equation*}
r_{1}=z_{1}^{k}+z_{2}^{k}+\cdots+z_{t-1}^{k}+r_{t}, \quad \text { where } \quad r_{t}<_{k} r_{1}^{\kappa^{t-1}}<_{k} m^{\kappa^{t}} \tag{5.4}
\end{equation*}
$$

and each $z_{i}$ is a nonnegative integer.
Let $x_{1}, \ldots, x_{R}$ represent integers of absolute value not exceeding $\sqrt{l}$, whose precise values will be chosen below. By Lemma 5.3, we can write

$$
\begin{align*}
M l^{k} & =x_{1}^{2 k}+\sum_{h=1}^{Q} u_{h}^{2 k} \\
M l^{k} & =x_{2}^{2 k}+\sum_{h=Q+1}^{2 Q} u_{h}^{2 k} \\
& \vdots  \tag{5.5}\\
M l^{k} & =x_{R}^{2 k}+\sum_{h=(R-1) Q+1}^{R Q} u_{h}^{2 k}
\end{align*}
$$

for certain integers $u_{1}, \ldots, u_{R Q}$. Adding equations (5.4) and (5.5), we find that

$$
\begin{aligned}
m & =R M l^{k}+r_{1} \\
& =\sum_{j=1}^{R} x_{j}^{2 k}+\sum_{h=1}^{Q R} u_{h}^{2 k}+z_{1}^{k}+z_{2}^{k}+\cdots+z_{t-1}^{k}+r_{t}
\end{aligned}
$$

We appear to have made some progress towards our goal, seeing as we have expressed $m$ as a sum of $R+Q R+t-1$ nonnegative $k$ th powers, up to a small remainder $r_{t}$. In particular, it would be an easy task to complete the proof if we knew that $r_{t}$ was the sum of a bounded number of $k$ th powers; however, this is not at all obvious.

To circumvent this difficulty we make a judicious choice of the numbers $x_{j}$. In the notation of Lemma 5.4, we set $x_{j}:=n+a_{j}^{\prime}$ for all $1 \leq j \leq R$, where $n$ is an integer which remains to be selected. Then $\sum_{j=1}^{R} x_{j}^{2 k}=$
$\sum_{i=1}^{R} y_{i}^{2 k}-(A n+B)$, where each $y_{i}:=n+a_{i}$. Hence

$$
\begin{equation*}
m=\sum_{i=1}^{R} y_{i}^{2 k}+\sum_{h=1}^{Q R} u_{h}^{2 k}+z_{1}^{k}+z_{2}^{k}+\cdots+z_{t-1}^{k}+r_{t}-(A n+B) . \tag{5.6}
\end{equation*}
$$

We now choose $n$ so that

$$
0 \leq r_{t}-(A n+B)<A, \quad \text { which amounts to setting } \quad n:=\left\lfloor\frac{r_{t}-B}{A}\right\rfloor .
$$

To see that we are permitted to choose $n$ in this way, we must check that each $x_{j}=n+a_{j}^{\prime}$ is at most $\sqrt{l}$ in absolute value. But by (5.4),

$$
x_{j}=\left\lfloor\frac{r_{t}-B}{A}\right\rfloor+a_{j}^{\prime} \ll k r_{t}+1 \ll_{k} m^{\kappa^{t}},
$$

while by (5.3), $\sqrt{l} \gg_{k} m^{\frac{1}{2 k}}$. Since $\kappa^{t}<\frac{1}{2 k}$, each $\left|x_{j}\right|$ is smaller than $\sqrt{l}$ if $m$ is sufficiently large (as we are assuming).

Since $0 \leq r_{t}-(A n+B)<A$, the integer $r_{t}-(A n+B)$ is a sum of fewer than $A$ terms of the form $1^{k}$. So by (5.6),

$$
g(k ; m)<R+Q R+t-1+A=O_{k}(1) .
$$

This completes the proof of Theorem 5.1.

## 3. Producing the Hilbert-Dress identities

3.1. Prerequisites from convex analysis. The proof of Theorem 5.2 given in this text assumes some familiarity with convex sets. Any number of sources would suffice for the the vocabulary and basic theory that we require; references below are to [Lay82].

Suppose that $V$ is a real vector space and that $S$ is a subset of $V$. We write conv $S$ for the convex hull of $S$. The following two results will be particularly important in what follows:
$\star$ Lemma 5.6 (Carathéodory's theorem). Let $V$ be a real vector space of dimension $n$. If $S$ is an arbitrary subset of $V$, then every element of conv $S$ can be written as a convex combination of at most $n+1$ elements of $S$. That is, if $\mathbf{v} \in \operatorname{conv} S$, then there is an $m \leq n$ for which $\mathbf{v}$ can be written in the form

$$
\begin{equation*}
\sum_{i=0}^{m} \alpha_{i} \mathbf{v}_{i}, \quad \text { where each } \mathbf{v}_{i} \in S, \quad \text { each } \alpha_{i} \geq 0, \quad \text { and } \quad \sum_{i=0}^{m} \alpha_{i}=1 \tag{5.7}
\end{equation*}
$$

Suppose, moreover, that with respect to some basis of $V$, not only the vector $\mathbf{v}$ but also all the vectors in $S$ have rational coordinates. Then we can choose a representation (5.7) of $\mathbf{v}$ where all the $\alpha_{i}$ are rational.

## Remarks.

1. We can always arrange to have $m=n$ in the representation (5.7). Indeed, if $m<n$, then (5.7) continues to hold with $m$ replaced by $n$ if we pad our representation by setting $\alpha_{i}:=0$ for $m<i \leq n$ and choose $\mathbf{v}_{i} \in S$ arbitrarily for these indices.
2. The second half of the lemma is often not stated explicitly in discussions of Carathéodory's theorem but is implicit in the usual proof of that result (see, e.g., [Lay82, Theorem 2.23]). Indeed, suppose that $\mathbf{v}$ and all the vectors in $S$ have rational coordinates, and write $\mathbf{v}$ in the form (5.7) with $m$ as small as possible. The minimality of $m$ forces $\mathbf{v}_{0}, \ldots, \mathbf{v}_{m}$ to be affinely independent (in the sense of $[\mathbf{L a y 8 2}$, Definition 2.17]). It follows that the real numbers $\alpha_{0}, \ldots, \alpha_{m}$ appearing in our representation (5.7) are the unique real numbers satisfying

$$
\mathbf{v}=\sum_{i=0}^{m} \alpha_{i} \mathbf{v}_{i}, \quad \text { where } \quad \sum_{i=0}^{m} \alpha_{i}=1
$$

By hypothesis, (5.8) defines a system of linear equations in the $\alpha_{i}$ with rational coefficients, so by Gaussian elimination its unique solution $\alpha_{0}, \ldots, \alpha_{m}$ must consist of rational numbers.

Lemma 5.7. Let $V$ be an n-dimensional real vector space and let $S$ be a convex subset of $V$. If the vector $\mathbf{v} \in V$ does not belong to the relative interior of $S$, then one can pass an ( $n-1$ )-dimensional hyperplane $H$ through $\mathbf{v}$ so that $S$ is contained entirely in one of the closed half-spaces determined by $H$.

Proof (sketch). By [Lay82, Corollary 4.6], there is an ( $n-1$ )-dimensional hyperplane $H$ through $\mathbf{v}$ with the relative interior of $S$ entirely contained in one of the open half-spaces determined by $H$. So the closure of the relative interior of $S$, which coincides with the closure of $S$ (cf. [Lay82, Exercise $2.13]$ ), belongs to the corresponding closed half-space.
3.2. Proof of Lemma 5.2. Let $V$ be the space of homogeneous polynomials of degree $2 k$ belonging to $\mathbf{R}\left[X_{1}, \ldots, X_{5}\right]$. Then $V$ is a real vector space of dimension $N:=\binom{2 k+4}{4}$, with a basis given by (an arbitrary ordering of) the monomials

$$
\begin{equation*}
X_{1}^{e_{1}} X_{2}^{e_{2}} X_{3}^{e_{3}} X_{4}^{e_{4}} X_{5}^{e_{5}}, \quad \text { where each } \quad e_{i} \geq 0 \quad \text { and } \quad \sum_{i=1}^{5} e_{i}=2 k \tag{5.9}
\end{equation*}
$$

We put an inner product on $V$ by using (5.9) to identify $V$ with $\mathbf{R}^{N}$ and then importing the standard dot product on $\mathbf{R}^{N}$. If $\left(\alpha_{1}, \ldots, \alpha_{5}\right) \in \mathbf{R}^{5}$, we
put

$$
\mathbf{v}_{\left(\alpha_{1}, \ldots, \alpha_{5}\right)}:=\left(\alpha_{1} X_{1}+\cdots+\alpha_{5} X_{5}\right)^{2 k} \in V
$$

Let $\mathbf{B}$ be the 5 -dimensional unit ball $\left\{\left(\alpha_{1}, \ldots, \alpha_{5}\right) \in \mathbf{R}^{5}: \alpha_{1}^{2}+\cdots+\alpha_{5}^{2} \leq 1\right\}$. Define a subset $S_{\text {B }}$ of $V$ by

$$
S_{\mathbf{B}}:=\left\{\mathbf{v}_{\left(\alpha_{1}, \ldots, \alpha_{5}\right)}:\left(\alpha_{1}, \ldots, \alpha_{5}\right) \in \mathbf{B}\right\}
$$

Lemma 5.8. Let $c=c(k)$ be the positive real number given by

$$
c:=\frac{\int_{\mathbf{B}} \beta_{1}^{2 k} d \beta_{1} d \beta_{2} \cdots d \beta_{5}}{\int_{\mathbf{B}} d \beta_{1} d \beta_{2} \cdots d \beta_{5}} .
$$

Then $f\left(X_{1}, \ldots, X_{5}\right):=c\left(X_{1}^{2}+\cdots+X_{5}^{2}\right)^{k}$ belongs to the relative interior of the convex hull of $S_{\mathbf{B}}$.

Proof. The proof proceeds in two stages. First we show that if we put

$$
\begin{align*}
g\left(X_{1}, \ldots, X_{5}\right): & =\frac{\int_{\mathbf{B}} \mathbf{v}_{\left(\alpha_{1}, \ldots, \alpha_{5}\right)} d \alpha_{1} \cdots d \alpha_{5}}{\int_{\mathbf{B}} d \alpha_{1} \cdots d \alpha_{5}} \\
& =\frac{\int_{\mathbf{B}}\left(\alpha_{1} X_{1}+\cdots+\alpha_{5} X_{5}\right)^{2 k} d \alpha_{1} \cdots d \alpha_{5}}{\int_{\mathbf{B}} d \alpha_{1} \cdots d \alpha_{5}} \tag{5.10}
\end{align*}
$$

then $f=g$ in $\mathbf{R}\left[X_{1}, \ldots, X_{5}\right]$. Then we show that $g$ belongs to the relative interior of the convex hull of $S_{\mathbf{B}}$.

Since two multivariate polynomials with real coefficients are equal if they agree for every assignment of the variables, to show that $f=g$ it is enough to prove that

$$
f\left(x_{1}, \ldots, x_{5}\right)=g\left(x_{1}, \ldots, x_{5}\right)
$$

for all real numbers $x_{1}, \ldots, x_{5}$. If all of the $x_{i}$ vanish, then $f=g=0$. Otherwise we perform the following change of variables in (5.10): Let

$$
\begin{aligned}
& \beta_{1}=\lambda_{11} \alpha_{1}+\lambda_{12} \alpha_{2}+\cdots+\lambda_{15} \alpha_{5}, \\
& \vdots \quad \vdots \quad \vdots \quad \ddots \quad \vdots \\
& \beta_{5}=\lambda_{51} \alpha_{1}+\lambda_{52} \alpha_{2}+\cdots+\lambda_{55} \alpha_{5},
\end{aligned}
$$

where

$$
\lambda_{i 1}:=\frac{x_{i}}{\sqrt{x_{1}^{2}+x_{2}^{2}+\cdots+x_{5}^{2}}}
$$

and the remaining $\lambda_{i j}$ are chosen so that the resulting matrix $\left(\lambda_{i j}\right)$ is orthogonal. The orthogonality of the matrix ensures that $\mathbf{B}$ is taken to itself by this linear transformation, and we find that

$$
\begin{aligned}
g\left(x_{1}, \ldots, x_{5}\right) & =\left(\frac{\int_{\beta_{1}^{2}+\cdots+\beta_{5}^{2} \leq 1} \beta_{1}^{2 k} d \beta_{1} d \beta_{2} \cdots d \beta_{5}}{\int_{\beta_{1}^{2}+\cdots+\beta_{5}^{2} \leq 1} d \beta_{1} d \beta_{2} \cdots d \beta_{5}}\right)\left(x_{1}^{2}+\cdots+x_{5}^{2}\right)^{k} \\
& =f\left(x_{1}, \ldots, x_{5}\right) .
\end{aligned}
$$

Thus $f=g$ for all choices of the $x_{i}$, and so $f=g$ in $\mathbf{R}\left[X_{1}, \ldots, X_{5}\right]$.
Now we take up the problem of showing that $g$ belongs to the relative interior of the convex hull of $S_{\mathbf{B}}$. Let $W$ be the affine hull of conv $S_{\mathrm{B}}$ (or equivalently, of $S_{\mathbf{B}}$ ). Since $S_{\mathbf{B}}$ contains the zero vector, $W$ coincides with the subspace of $V$ generated by $S_{\mathbf{B}}$. We would like to apply Lemma 5.7 (with $\mathbf{v}=g$ and $V=W$ ), but the way we have set things up, it is necessary to first verify that $g$ belongs to $W$. But this is easy: Indeed, if $\mathbf{v}$ is any vector from the orthogonal complement $W^{\perp}$ of $W$, then $\mathbf{v} \cdot g=0$, since $\mathbf{v} \cdot \mathbf{v}_{\left(\alpha_{1}, \ldots, \alpha_{5}\right)}=0$ for all $\left(\alpha_{1}, \ldots, \alpha_{5}\right) \in \mathbf{B}$. So $g \in\left(W^{\perp}\right)^{\perp}=W$.

So by Lemma 5.7, if $g$ does not belong to the relative interior of conv $S_{\mathrm{B}}$, then there is a hyperplane $H$ in $W$ passing through $g$ with conv $S_{\mathbf{B}}$ entirely contained in one of the closed half-spaces determined by $H$. Such a hyperplane corresponds to a nonzero $\mathbf{w} \in W$ with the property that

$$
\mathbf{w} \cdot \mathbf{v}_{\left(\alpha_{1}, \ldots, \alpha_{5}\right)} \geq \mathbf{w} \cdot g
$$

for all $\left(\alpha_{1}, \ldots, \alpha_{5}\right) \in \mathbf{B}$. But then

$$
\mathbf{w} \cdot g=\frac{\int_{\mathbf{B}}\left(\mathbf{w} \cdot \mathbf{v}_{\left(\alpha_{1}, \ldots, \alpha_{5}\right)}\right) d \alpha_{1} \cdots d \alpha_{5}}{\int_{\mathbf{B}} d \alpha_{1} \cdots d \alpha_{5}} \geq \frac{\int_{\mathbf{B}}(\mathbf{w} \cdot g) d \alpha_{1} \cdots d \alpha_{5}}{\int_{\mathbf{B}} d \alpha_{1} \cdots d \alpha_{5}}=\mathbf{w} \cdot g
$$

which forces us to have

$$
\mathbf{w} \cdot \mathbf{v}_{\left(\alpha_{1}, \ldots, \alpha_{5}\right)}=\mathbf{w} \cdot g
$$

for all $\left(\alpha_{1}, \ldots, \alpha_{5}\right) \in \mathbf{B}$. Since $\mathbf{B}$ contains $(0,0,0,0,0)$, this implies that $\mathbf{w}$ is orthogonal to $\mathbf{v}_{\left(\alpha_{1}, \ldots, \alpha_{5}\right)}$ for all $\left(\alpha_{1}, \ldots, \alpha_{5}\right) \in \mathbf{B}$. Thus $S_{\mathbf{B}}$ is entirely contained within a proper subspace of $W$ (viz. the hyperplane orthogonal to w), contrary to the choice of $W$.

Lemma 5.8 is enough to prove the existence of an identity of the form (5.1) but where $M$, the $m_{i}$, and the $a_{i j}$ are real numbers (and not necessarily integers). In order to obtain Theorem 5.2 as stated, it is expedient to introduce the following relatives of $S$ :

$$
S_{\mathbf{R}^{5}}:=\left\{\mathbf{v}_{\left(\alpha_{1}, \ldots, \alpha_{5}\right)}: \alpha_{i} \in \mathbf{R}\right\} \quad \text { and } \quad S_{\mathbf{Q}^{5}}:=\left\{\mathbf{v}_{\left(\alpha_{1}, \ldots, \alpha_{5}\right)}: \alpha_{i} \in \mathbf{Q}\right\} .
$$

Theorem 5.2 will follow once we know that the $f$ of Lemma 5.8 belongs not only to the relative interior of conv $S_{\mathbf{B}}$ but also to the relative interior of conv $S_{\mathbf{Q}^{5}}$. This is a consequence of the next two lemmas:

Lemma 5.9. The sets $S_{\mathbf{B}}$ and $S_{\mathbf{R}^{5}}$ have the same affine hull.
Proof. Since $\mathbf{0} \in S_{\mathbf{B}} \cap S_{\mathbf{R}^{5}}$, this amounts to checking that $S_{\mathbf{B}}$ and $S_{\mathbf{R}^{5}}$ generate the same subspace of $V$. But this is clear, since $S_{\mathbf{R}^{5}}$ is a union of dilations of $S_{\mathbf{B}}$.

Since $S_{\mathbf{B}}$ is contained within $S_{\mathbf{R}^{5}}$, its convex hull conv $S_{\mathbf{B}}$ is contained within conv $S_{\mathbf{R}^{5}}$. So by Lemma 5.9, the relative interior of $\operatorname{conv} S_{\mathbf{B}}$ is contained in the relative interior of conv $S_{\mathbf{R}^{5}}$. Since $f$ belongs to the relative interior of conv $S_{\mathbf{B}}$ by Lemma 5.8, we will have that $f$ belongs to the relative interior of conv $S_{\mathbf{Q}^{5}}$ once we establish the following lemma:
Lemma 5.10. The relative interior of conv $S_{\mathbf{R}^{5}}$ is contained within the relative interior of conv $S_{\mathbf{Q}^{5}}$.

Proof. We start by observing that, using an overline to denote the closure operator,

$$
\begin{equation*}
\operatorname{conv} S_{\mathbf{R}^{5}} \subset \overline{\operatorname{conv} S_{\mathbf{Q}^{5}}} \tag{5.11}
\end{equation*}
$$

Indeed, suppose that $\mathbf{v}$ belongs to the convex hull of $S_{\mathbf{R}^{5}}$ and write $\mathbf{v}$ as a convex combination of vectors $\mathbf{v}_{i} \in S_{\mathbf{R}^{5}}$. We can approximate these $\mathbf{v}_{i}$ arbitrarily closely by elements of $S_{\mathbf{Q}^{5}}$, and so we can approximate $\mathbf{v}$ arbitrarily closely by elements of conv $S_{\mathbf{Q}^{5}}$. This proves (5.11).

Consequently, the affine hull of conv $S_{\mathbf{R}^{5}}$ is contained within the affine hull of $\overline{\text { conv } S_{\mathbf{Q}^{5}}}$. On the other hand, the affine hull of conv $S_{\mathbf{R}^{5}}$ coincides with the affine hull of $S_{\mathbf{R}^{5}}$ while the affine hull of $\overline{\operatorname{conv} S_{\mathbf{Q}^{5}}}$ coincides with the affine hull of $S_{\mathbf{Q}^{5}}$. Since $S_{\mathbf{Q}^{5}}$ is contained in $S_{\mathbf{R}^{5}}$, we conclude that both sides of (5.11) have the same affine hull.

It now follows from (5.11) that the relative interior of conv $S_{\mathbf{R}^{5}}$ is contained in the relative interior of $\overline{\operatorname{conv} S_{\mathbf{Q}^{5}}}$. To complete the proof of the lemma we need only recall that a convex set and its closure always have the same relative interior (cf. [Lay82, Exercise 2.14]).

Proof of Theorem 5.2. Since $\mathbf{0} \in S_{\mathbf{Q}^{5}}$, the affine hull of $S_{\mathbf{Q}^{5}}$ coincides with the subspace generated by $S_{\mathbf{Q}^{5}}$. Since $X_{5}^{2 k}$ belongs to this subspace (because $\left.X_{5}^{2 k}=\mathbf{v}_{(0,0,0,0,1)} \in S_{\mathbf{Q}^{5}}\right)$ and $f\left(X_{1}, \ldots, X_{5}\right)$ is in the relative interior of conv $S_{\mathbf{Q}^{5}}$, it follows that for all sufficiently small $\mu>0$,

$$
f\left(X_{1}, \ldots, X_{5}\right)-\mu X_{5}^{2 k}=c\left(X_{1}^{2}+X_{2}^{2}+\cdots+X_{5}^{2}\right)^{k}-\mu X_{5}^{2 k} \in \operatorname{conv} S_{\mathbf{Q}^{5}}
$$

Moreover, since $\mathbf{0} \in \operatorname{conv} S_{\mathbf{Q}^{5}}$, for each $0 \leq \lambda \leq 1$,

$$
\lambda c\left(X_{1}^{2}+X_{2}^{2}+\cdots+X_{5}^{2}\right)^{k}-\lambda \mu X_{5}^{2 k} \in \operatorname{conv} S_{\mathbf{Q}^{5}} .
$$

We choose $\lambda>0$ and $\mu$ here so that both $\lambda c$ and $\lambda \mu$ are rational. Applying Carathéodory's Theorem (Theorem 5.6), we can write

$$
\lambda c\left(X_{1}^{2}+X_{2}^{2}+\cdots+X_{5}^{2}\right)^{k}=\lambda \mu X_{5}^{2 k}+\sum_{i=0}^{N} r_{i}\left(b_{i 1} X_{1}+\cdots+b_{i 5} X_{5}\right)^{2 k}
$$

where each $r_{i} \geq 0$ is rational, $\sum_{i=0}^{N} r_{i}=1$, and each $b_{i j}$ is rational. The Hilbert-Dress identity follows upon clearing all the denominators.

## Notes

Our proof of Theorem 5.1 is a pure existence argument; it shows that $g(k)$ is bounded above but does not give any finite procedure for determining an upper bound. This is because our proof of Lemma 5.2, due essentially to Hilbert [Hil09] and Schmidt [Sch13], yields no information on the size of the coefficients in (5.1). An alternative method for proving identities like (5.1) was given by Hausdorff [Hau09]. This allowed Rieger [Rie53a, Rie53b], in his doctoral dissertation, to obtain explicit upper bounds on $g(k)$. Specifically, he proved that for each $k$,

$$
g(k)<(2 k+1)^{260(k+3)^{3 k+8}} .
$$

He later announced in [Rie56] the improved bound

$$
g(k)<(2 k+1)^{260(k+1)^{8}} .
$$

If instead of following Hilbert's original proof, one uses Rieger's method in combination with Dress's argument, then one finds that

$$
g(k)<(2 k+1)^{2000 k^{5}}
$$

(see [Pol09]). This appears to be the best known general upper bound on $g(k)$ so far obtained by elementary methods, although as we shall see shortly, it is quite far from the truth.

Around 1772, J. A. Euler observed that $g(k ; n)=2^{k}+\left\lfloor(3 / 2)^{k}\right\rfloor-2$ when $n:=2^{k}\left\lfloor(3 / 2)^{k}\right\rfloor-1$. (The reader should attempt to verify this for herself; the essential observation is that $n<3^{k}$, so that only $0^{k}, 1^{k}$, and $2^{k}$ can enter into a representation of $n$ as a sum of $k$ th powers.) Thus

$$
\begin{equation*}
g(k) \geq 2^{k}+\left\lfloor(3 / 2)^{k}\right\rfloor-2 \tag{5.12}
\end{equation*}
$$

In particular, $g(2) \geq 4, g(3) \geq 9, g(4) \geq 19, g(5) \geq 53$, etc. Remarkably, it turns out that the easy lower bound (5.12) is almost always sharp. More precisely, we have the following statement (which combines results of Dickson, Pillai, Rubugunday, Niven, Chen, Balasubramanian, Deshouillers, and Dress): Write $\{x\}$ for the fractional part $x-\lfloor x\rfloor$ of the real number $x$. If

$$
\begin{equation*}
2^{k}\left\{(3 / 2)^{k}\right\}+\left\lfloor(3 / 2)^{k}\right\rfloor \leq 2^{k}, \tag{5.13}
\end{equation*}
$$

then equality holds in (5.12). If (5.13) fails, define

$$
N(k):=\left\lfloor(3 / 2)^{k}\right\rfloor \cdot\left\lfloor(4 / 3)^{k}\right\rfloor+\left\lfloor(3 / 2)^{k}\right\rfloor+\left\lfloor(4 / 3)^{k}\right\rfloor ;
$$

then

$$
g(k):= \begin{cases}\left\lfloor(3 / 2)^{k}\right\rfloor+\left\lfloor(4 / 3)^{k}\right\rfloor+2^{k}-3 & \text { if } 2^{k}<N(k), \\ \left\lfloor(3 / 2)^{k}\right\rfloor+\left\lfloor(4 / 3)^{k}\right\rfloor+2^{k}-2 & \text { if } 2^{k}=N(k) .\end{cases}
$$

The inequality (5.13) holds for all $k \leq 471,600,000[\mathbf{K W 9 0}]$, and it seems reasonable to conjecture that it always holds. In any event, Mahler [Mah57] has shown that (5.13) fails for at most finitely many $k$.

Much of the modern research on Waring's problem focuses not on $g(k)$, but on the quantity $G(k)$, defined as the smallest number of $k$ th powers needed to additively represent all sufficiently large numbers. (Thus $G(k)=$ $\lim \sup _{n \rightarrow \infty} g(k ; n)$.) For $k=2$, we have $g(2)=G(2)=4$, since no number from the residue class $7 \bmod 8$ is a sum of three squares. But for $k>2$, it is known that $G(k)<g(k)$. In fact, for large $k, G(k)$ is considerably smaller than $g(k)$; in sharp contrast with (5.12), Wooley [Woo95] has proved that

$$
G(k) \leq k \log k+k \log \log k+2 k+O(k \log \log k / \log k) .
$$

The precise determination of $G(k)$ is a very difficult problem which has been solved only for $k=2$ and $k=4$ (see [Dav39]).

The results of the last two paragraphs rely on the Hardy-Littlewood circle method for their proofs. For a gentle introduction to this method in the context of Waring's problem, see [Nat96, Chapters 4 and 5]. For further discussion of Waring's problem, see [HW08, Chapter XXI] and the surveys of Ellison [E1171] and Vaughan \& Wooley [VW02].

## Sieve Methods

Brun's [sieve] method ... is perhaps our most powerful elementary tool in number theory. - P. Erdős [Erd65]

## 1. Introduction

1.1. The sieve of Eratosthenes. Granville has pointed out [Gra95] that ancient Greek mathematics produced two results in prime number theory that have proved of first importance in subsequent thought. The first is Euclid's proof of the infinitude of the primes, which was discussed in Chapter 1. The second is the sieve of Eratosthenes.

Eratosthenes' method allows one to determine the primes not exceeding $x$ assuming only knowledge of the primes not exceeding $\sqrt{x}$. In this procedure one begins with a list of all positive integers in the interval $[2, x]$. For each prime $p \leq \sqrt{x}$, we cross out all the multiples of $p$ on the list; the numbers remaining are exactly the primes in the interval $(\sqrt{x}, x]$. We illustrate this with $x=30$, sieving by the primes 2,3 , and 5 :

$$
\begin{array}{cccccccccc} 
& 2 & 3 & 4 & 5 & 8 & 7 & 8 & 9 & 10 \\
11 & 12 & 13 & 14 & 15 & 16 & 17 & 18 & 19 & 20 \\
21 & 22 & 23 & 24 & 25 & 26 & 27 & 28 & 29 & 30
\end{array}
$$

This procedure is remarkable not only insofar as it gives a fast algorithm for listing primes, but also in that it suggests the useful viewpoint of the primes as the integers surviving a "sieving process".
1.2. Legendre's formula. Let us attempt to count how many integers remain after Eratosthenes' sieving procedure is carried out. More generally, let us count the number of positive integers up to $x$ remaining after the
deletion (or "sifting out") of the multiples of all primes not exceeding $z$, where $z$ is a parameter at our disposal (in Eratosthenes' sieve, $z=\sqrt{x}$ ). We use $\pi(x, z)$ to denote this quantity, i.e.,

$$
\pi(x, z):=\#\{n \leq x: p \mid n \Rightarrow p>z\} .
$$

Then for every $z>0$,

$$
\pi(x) \leq z+\pi(x, z)
$$

and

$$
\pi\left(x, x^{1 / 2}\right)=\pi(x)-\pi(\sqrt{x})+1 .
$$

Our estimate of $\pi(x, z)$ proceeds in several stages. We begin with the total number $\lfloor x\rfloor$ of positive integers not exceeding $x$, and then for each prime $p_{1} \leq z$ we subtract the number of multiples of $p_{1}$ :

$$
\lfloor x\rfloor-\sum_{p_{1} \leq z}\left\lfloor\frac{x}{p_{1}}\right\rfloor .
$$

This counts correctly those $n$ with at most one prime divisor $p \leq z$, but those $n$ with two or more prime factors $p \leq z$ have been subtracted off twice. Hence, we add these back in, to obtain our next approximation,

$$
\lfloor x\rfloor-\sum_{p_{1} \leq z}\left\lfloor\frac{x}{p_{1}}\right\rfloor+\sum_{p_{1}<p_{2} \leq z}\left\lfloor\frac{x}{p_{1} p_{2}}\right\rfloor .
$$

But now those integers divisible by three primes $p \leq z$ have been added back in too many times; for instance, if $n$ has exactly three prime divisors not exceeding $z$, it is counted with weight $1-3+3>0$. Thus we should subtract a term corresponding to the integers divisible by three primes $p \leq z$; we would then find ourselves needing to add a term corresponding to integers divisible by four such $p$, etc. Continuing in this manner, we are led to the formula

$$
\begin{equation*}
\pi(x, z)=\lfloor x\rfloor-\sum_{p_{1} \leq z}\left\lfloor\frac{x}{p_{1}}\right\rfloor+\cdots+(-1)^{r} \sum_{p_{1}<\cdots<p_{r} \leq z}\left\lfloor\frac{x}{p_{1} \cdots p_{r}}\right\rfloor, \tag{6.1}
\end{equation*}
$$

where $r=\pi(z)$. If we set

$$
P:=\prod_{p \leq z} p,
$$

we can put (6.1) in the alternative form

$$
\begin{equation*}
\pi(x, z)=\sum_{d \mid P} \mu(d)\left\lfloor\frac{x}{d}\right\rfloor . \tag{6.2}
\end{equation*}
$$

This reasoning, due to Legendre, can be tightened into a proof of (6.1). For the time being we assume (6.1), postponing a rigorous justification to §3, where we will establish a more general result.
1.3. Consequences. We now have an exact formula for $\pi(x, z)$. Unfortunately this exact formula is a bit unsatisfying because it leaves the most natural question unanswered: How large is $\pi(x, z)$ ?

What does our formula (6.2) have to say about this question? Sums involving the greatest-integer function are generally difficult to work with, so we drop the greatest integer signs in (6.2) and transfer the incurred error to a separate sum:

$$
\pi(x, z)=x \sum_{d \mid P} \frac{\mu(d)}{d}+\sum_{d \mid P} \mu(d)\left(\left\lfloor\frac{x}{d}\right\rfloor-\frac{x}{d}\right) .
$$

The first sum can be written as the product $\prod_{p \leq z}(1-1 / p)$. The second sum (which we view as the error term) is bounded in absolute value by $2^{\pi(z)}$, since there are $2^{\pi(z)}$ divisors $d$ of $P$, and for each of these the corresponding summand has absolute value at most 1 . Thus

$$
\begin{equation*}
\pi(x, z)=x \prod_{p \leq z}\left(1-\frac{1}{p}\right)+O\left(2^{\pi(z)}\right) . \tag{6.3}
\end{equation*}
$$

How useful is this estimate? Suppose first that $z$ is fixed while $x$ is tending to infinity; then the error term in (6.3) is $O_{z}(1)$ and we obtain the asymptotic formula $\pi(x, z) \sim x \prod_{p \leq z}(1-1 / p)$. The same asymptotic estimate holds if $z$ is not fixed, but instead tends to infinity with $x$ sufficiently slowly.

Whenever $z=z(x) \rightarrow \infty$, Mertens' theorem implies that

$$
\begin{equation*}
x \prod_{p \leq z}\left(1-\frac{1}{p}\right) \sim e^{-\gamma} \frac{x}{\log z} \quad(x \rightarrow \infty) \tag{6.4}
\end{equation*}
$$

If $z=z(x)$ also satisfies $z \leq \log x$ once $x$ is sufficiently large, then the $O$-term in (6.3) is $\ll 2^{z} \leq x^{\log 2}$, which is of smaller order than $x / \log z$. Consequently, $\pi(x, z) \sim e^{-\gamma} x / \log z$. Taking $z=\log x$, we obtain that

$$
\begin{equation*}
\pi(x) \leq \pi(x, \log x)+\log x \leq\left(e^{-\gamma}+o(1)\right) \frac{x}{\log \log x}, \tag{6.5}
\end{equation*}
$$

which provides another proof that the set of primes has density zero.
We have yet to treat the case corresponding to Eratosthenes' sieve, that of $z=\sqrt{x}$. In this case the "main term" in (6.3) is

$$
\begin{equation*}
x \prod_{p \leq x^{1 / 2}}\left(1-\frac{1}{p}\right) \sim 2 e^{-\gamma} \frac{x}{\log x}=(1.229 \ldots) \frac{x}{\log x} . \tag{6.6}
\end{equation*}
$$

Unfortunately our bound of $2^{\pi(\sqrt{x})}$ for the "error term" dwarfs the value of this main term. (For example, by Chebyshev's results, $2^{\pi(\sqrt{x})}>2^{\sqrt[3]{x}}$ for large $x$, and $2 \sqrt[3]{x}^{x}$ grows faster than any fixed power of $x$.) So (6.3) does
not give us the asymptotic formula $\pi\left(x, x^{1 / 2}\right) \sim x \prod_{p \leq x^{1 / 2}}(1-1 / p)$. And in fact, by the prime number theorem,

$$
\begin{equation*}
\pi\left(x, x^{1 / 2}\right)=\pi(x)-\pi(\sqrt{x})+1 \sim x / \log x, \tag{6.7}
\end{equation*}
$$

so that it is not even true that $\pi\left(x, x^{1 / 2}\right) \sim x \prod_{p \leq x^{1 / 2}}(1-1 / p)$. This points to a limitation of our method for approximating $\pi(x, z)$; in $\S 1.5$ we will discuss to what extent difficulties of this sort can be overcome.
1.4. General sieving situations. The problem treated in the last section is of the following form: Given a finite sequence ${ }^{1}$ of integers $\mathcal{A}$ and a finite set of primes $\mathcal{P}$, estimate the number $S(\mathcal{A}, \mathcal{P})$ of terms of $\mathcal{A}$ divisible by no prime $p \in \mathcal{P}$. For example, if

$$
\begin{equation*}
\mathcal{A}:=\{n \leq x\} \quad \text { and } \quad \mathcal{P}:=\{p \leq z\}, \tag{6.8}
\end{equation*}
$$

then $S(\mathcal{A}, \mathcal{P})$ is what we have been calling $\pi(x, z)$.
Many problems in number theory fit into this framework. For example, suppose $x, z>0$. Set

$$
\begin{equation*}
\mathcal{A}:=\{n(n+2): n \leq x\}, \quad \mathcal{P}:=\{p \leq z\} . \tag{6.9}
\end{equation*}
$$

If both $n$ and $n+2$ are prime, then either $n \leq z$ or both $n$ and $n+2$ have only prime factors exceeding $z$. Consequently,

$$
\pi_{2}(x) \leq S(\mathcal{A}, \mathcal{P})+z .
$$

Moreover, $n$ and $n+2$ are both prime if all of their prime factors exceed $\sqrt{x+2}$. So if we take $z=\sqrt{x+2}$, then

$$
\begin{equation*}
0 \leq \pi_{2}(x)-S(\mathcal{A}, \mathcal{P}) \leq z \tag{6.10}
\end{equation*}
$$

Estimates for $S(\mathcal{A}, \mathcal{P})$ are thus intimately connected with the quantitative version of the twin prime conjecture introduced in Chapter $3, \S 5$.

In order to prove any general theorems on the size of $S(\mathcal{A}, \mathcal{P})$, it is necessary to make some further assumptions. We will assume that $\mathcal{A}$ has "approximate" length $X$ and that divisibility by distinct primes $p \in \mathcal{P}$ constitute "approximately" independent events, each occurring with "approximate" probability $\alpha(p)$. (All of this will be made precise in §2.) In this case, it is natural to expect that

$$
\begin{equation*}
S(\mathcal{A}, \mathcal{P}) \approx X \prod_{p \in \mathcal{P}}(1-\alpha(p)) \tag{6.11}
\end{equation*}
$$

From our perspective in this chapter, the goal of sieve theory is to quantify and then to justify such approximations, in as wide a range of circumstances as possible.

[^3]In the classical situation described by (6.8), it is reasonable to approximate the count of natural numbers $n \leq x$ by $x$ and the probability that such an integer is divisible by $p$ by $1 / p$. Then (6.11) is the guess that $\pi(x, z) \approx x \prod_{p \leq z}(1-1 / p)$. We have seen that when $z$ is constant or slowgrowing, this approximation holds as an asymptotic formula, but that for $z=\sqrt{x}$ (the case originally of interest), the approximation is off by a constant factor. Nevertheless, (6.11) is still correct if read as the assertion that both sides have the same order of magnitude.

For another example, consider the situation described by (6.9). Again the length of $\mathcal{A}$ is approximately $x$. The probability that a term of $\mathcal{A}$ is divisible by the prime $p$ is approximately $\nu(p) / p$, where $\nu(2)=1$ and $\nu(p)=2$ for $p>2$. (So that $\nu(p)$ is the number of solutions to $n(n+2) \equiv$ $0(\bmod p)$.$) The prediction (6.11) is that$

$$
\begin{equation*}
S(\mathcal{A}, \mathcal{P}) \approx x \prod_{p \leq z}\left(1-\frac{\nu(p)}{p}\right) \tag{6.12}
\end{equation*}
$$

If $z=z(x) \rightarrow \infty$ as $x \rightarrow \infty$, it is an easy deduction from Mertens' theorem (given in $\S 3.2$ below) that

$$
\begin{equation*}
x \prod_{p \leq z}\left(1-\frac{\nu(p)}{p}\right) \sim 2 C_{2} e^{-2 \gamma} \frac{x}{(\log z)^{2}} \tag{6.13}
\end{equation*}
$$

where $C_{2}=\prod_{p>2}\left(1-(p-1)^{-2}\right)$ is the twin prime constant. Arguing as in $\S 1.3$, one can show that $S(\mathcal{A}, \mathcal{P})$ is asymptotic to the right-hand side of (6.13) when $z$ is quite small (say $z \leq \frac{1}{2} \log x$ ) and $x$ tends to infinity. Probably no method can establish the same if $z=\sqrt{x+2}$; indeed, referring back to (6.10), we see that this would contradict the quantitative form of the twin prime conjecture (Conjecture 3.18). Note that even if $z=\sqrt{x+2}$, we still expect that the right-hand side of (6.13) has the same order of magnitude as $\pi_{2}(x)$; it is only off from what is conjecturally correct by a factor of $\left(2 e^{-\gamma}\right)^{2}$; cf. Exercise 28.
1.5. Legendre, Brun, and Hooley; oh my! We have already stated that the goal of sieve theory, for us, is to quantify and to justify estimates of the form

$$
S(\mathcal{A}, \mathcal{P}) \approx X \prod_{p \in \mathcal{P}}(1-\alpha(p))
$$

We can get a feel for the respective power of the three sieve methods considered in this chapter if we reflect on what they say about the particular estimate $\pi(x, z) \approx x \prod_{p \leq z}(1-1 / p)$ corresponding to our initial problem. As noted above, Legendre's method of successive approximation can be developed to show that this approximation is asymptotically correct when $z=\log x$. The first improvement on Legendre's methods, known as Brun's
pure sieve, shows that this remains true in a wider range: We need only assume that $z=z(x) \rightarrow \infty$ subject to the inequality $z(x) \leq x^{1 /(10 \log \log x)}$ (for large $x$ ). In particular, choosing $z$ as large as possible and referring to (6.4), we find that

$$
\begin{equation*}
\pi(x) \leq \pi(x, z)+z \ll \frac{x}{\log x} \log \log x \tag{6.14}
\end{equation*}
$$

which is considerably sharper than (6.5).
The final method to be developed in this chapter, known as the BrunHooley sieve, allows one to obtain upper and lower bounds for $\pi(x, z)$ when $z$ is as large as a (small) fixed power of $x$. From its upper bound we recover Chebyshev's estimate $\pi(x) \ll x / \log x$. (But one should take this with a grain of salt - in the derivation, we require the results of Mertens, which in turn rest on those of Chebyshev.) The lower bound aspect is also interesting, and allows one to deduce bounds of the shape $\pi\left(x, x^{1 / 1000}\right) \gg x / \log x$. Such a lower bound does not translate into a lower bound on $\pi(x)$; but because an integer up to $x$ all of whose prime factors exceed $x^{1 / 1000}$ can have at most 1000 prime factors, it does give us a lower bound on the number of 1000 -almost primes up to $x$. Here an $r$-almost prime is an integer with no more than $r$ prime divisors, counting multiplicity.

All of this might seem a bit silly because we have known the correct order of magnitude of $\pi(x)$ since Chapter 3. But the general sieve framework is rather flexible, and therein lies the potential of this approach. We have already seen that sieve methods can be adapted to yield information about the twin prime conjecture. Developing these ideas, Brun used his pure sieve to prove (in analogy with (6.14)) that

$$
\begin{equation*}
\pi_{2}(x) \ll \frac{x}{(\log x)^{2}}(\log \log x)^{2} . \tag{6.15}
\end{equation*}
$$

This is off by a factor of $(\log \log x)^{2}$ from the conjectured order of magnitude, but it still has profound implications. One consequence is that $\sum_{p} 1 / p$, restricted to primes $p$ which belong to a twin prime pair, is either a finite sum or a convergent infinite series.

Brun succeeded in removing the unwanted factor $(\log \log x)^{2}$ from (6.15) but required a rather complicated combinatorial apparatus to do so. We will reach the same goal by making use of simple ideas of Hooley. The same method will allow us to prove the following two deep theorems of Brun ([Bru20]; see [Wan84] for an English translation), approximations to the twin prime and Goldbach conjectures respectively:

- There are infinitely many pairs of 9 -almost primes $n, n+2$.
- Every large even integer $N$ is a sum of two 9 -almost primes.

In the next section we formally introduce some notions and notation arising in the general sieving situation. We then discuss the first sieve method, that of Eratosthenes-Legendre. This is just a general version of Legendre's method of successive approximation, seen above. After giving a few elementary applications, we turn to a discussion of Brun's pure sieve. This method gets its name from its origin in the purely combinatorial observation that the approximations in Legendre's method are alternately over and underestimates. Brun's pure sieve is much more powerful than Legendre's method, which we illustrate by proving the aforementioned theorem of Brun on the sum of the reciprocals of the twin primes. We then describe Hooley's elegant and surprisingly powerful "almost-pure" sieve, basing our treatment on Hooley's original article [Hoo94] and the exposition of Ford \& Halberstam [FH00]. We conclude the chapter with a striking application of sieve methods to the Goldbach problem, found by Schnirelmann.

## 2. The general sieve problem: Notation and preliminaries

Probability is not a notion of pure mathematics, but of philosophy or physics. - G. H. Hardy \& J. E. Littlewood [HL23]
The general sieve problem takes the following form: Given a finite sequence of integers $\mathcal{A}=\left\{a_{i}\right\}$ and a finite set of primes $\mathcal{P}$, estimate the quantity

$$
S(\mathcal{A}, \mathcal{P}):=\#\{a \in \mathcal{A}: \operatorname{gcd}(a, P)=1\}
$$

where $P:=\prod_{p \in \mathcal{P}} p$.
In many situations, the sieving set $\mathcal{P}$ is obtained by truncating an infinite set of primes at a point $z$. Consequently, it is expedient to allow the set $\mathcal{P}$ to be infinite and to introduce special notation indicating that we sieve only by those primes $p \in \mathcal{P}$ with $p \leq z$. We therefore define

$$
S(\mathcal{A}, \mathcal{P}, z):=\#\{a \in \mathcal{A}: \operatorname{gcd}(a, P(z))=1\}
$$

where

$$
P(z):=\prod_{\substack{p \in \mathcal{P} \\ p \leq z}} p .
$$

Hence $S(\mathcal{A}, \mathcal{P}, z)=S(\mathcal{A}, \mathcal{P} \cap[2, z])$.
We use the notation $A_{d}$ to denote the number of terms of $\mathcal{A}$ divisible by $d$, i.e.,

$$
A_{d}:=\#\{a \in \mathcal{A}: d \mid a\} .
$$

The letter $X$ denotes an approximation to the size of $\mathcal{A}$. We assume the existence of a multiplicative function $\alpha$ taking values in $[0,1]$ for which

$$
\begin{equation*}
A_{d}=X \alpha(d)+r(d) \tag{6.16}
\end{equation*}
$$

for each $d \mid P$ (or each $d \mid P(z)$, as the case may be). In practice, we choose $X$ and $\alpha$, and we define $r(d)$, for $d \mid P$, so that (6.16) holds.

## 3. The sieve of Eratosthenes-Legendre and its applications

3.1. The principle of inclusion-exclusion. Any rigorous study of sieve methods begins with the following fundamental result from enumerative combinatorics:

Theorem 6.1 (Principle of inclusion-exclusion). Let $X$ be a nonempty, finite set of $N$ objects, and let $P_{1}, \ldots, P_{r}$ be properties that elements of $X$ may have. For each subset $I \subset\{1,2, \ldots, r\}$, let $N(I)$ denote the number of elements of $X$ that have each of the properties indexed by the elements of $I$. Then with $N_{0}$ denoting the number of elements of $X$ with none of these properties, we have

$$
\begin{align*}
N_{0} & =\sum_{k=0}^{r}(-1)^{k} \sum_{\substack{I \subset\{1,2, \ldots, r\} \\
|I|=k}} N(I) \\
& =\sum_{I \subset\{1,2, \ldots, r\}}(-1)^{|I|} N(I) . \tag{6.17}
\end{align*}
$$

Proof. Suppose $x \in X$ has exactly $l$ of the properties $P_{1}, \ldots, P_{r}$. If $l=0$, then $x$ is counted only once in (6.17), in the term $N(\emptyset)$. On the other hand, if $1 \leq l \leq r$, then the number of $k$-element sets $I \subset\{1,2,3, \ldots, r\}$ for which $x$ is counted in $N(I)$ is exactly $\binom{l}{k}$, and the total weight with which $x$ is counted is

$$
\sum_{k=0}^{l}(-1)^{k}\binom{l}{k}=(1-1)^{l}=0
$$

by the binomial theorem.
3.2. A first sieve result. The principle of inclusion-exclusion can be applied immediately to the situation of $\S 2$ :

Theorem 6.2 (Sieve of Eratosthenes-Legendre).

$$
S(\mathcal{A}, \mathcal{P})=X \prod_{p \in \mathcal{P}}(1-\alpha(p))+\sum_{d \mid P} \mu(d) r(d) .
$$

Proof. Let $p_{1}, \ldots, p_{r}$ be a list of the primes in $\mathcal{P}$, and for each $1 \leq i \leq r$, let $P_{i}$ be the property of being divisible by $p_{i}$. For every $d \mid P$, there are
$X \alpha(d)+r(d)$ terms $a \in \mathcal{A}$ divisible by $d$. So by the principle of inclusionexclusion, the number of $a \in \mathcal{A}$ divisible by none of the primes of $\mathcal{P}$ is

$$
\begin{aligned}
& \sum_{k=0}^{r}(-1)^{k} \sum_{\substack{I \subset\{1,2, \ldots, r\} \\
|I|=k}} N(I)=\sum_{k=0}^{r}(-1)^{k} \sum_{\substack{d \mid P \\
\omega(d)=k}} A_{d} \\
& =\sum_{k=0}^{r} \sum_{\substack{d \mid P \\
\omega(d)=k}} \mu(d)(X \alpha(d)+r(d))=X \sum_{d \mid P} \mu(d) \alpha(d)+\sum_{d \mid P} \mu(d) r(d) \\
& =X \prod_{p \in \mathcal{P}}(1-\alpha(p))+\sum_{d \mid P} \mu(d) r(d) .
\end{aligned}
$$

Example. Let $\mathcal{A}=\{n \leq x\}$ and let $\mathcal{P}=\{p \leq z\}$. Then $S(\mathcal{A}, \mathcal{P})$ is what we referred to in the introduction as $\pi(x, z)$. For each $d$, we have $A_{d}=\lfloor x / d\rfloor$. So if we set $X=x$ and $\alpha(d)=1 / d$, and define $r(d)$ by (6.16), then $r(d)=-\{x / d\}$. In particular, $|r(d)| \leq 1$ for each $d$. So applying Theorem 6.2 with this choice of $X$ and $\alpha$, we recover the estimate (6.3), which was derived in a nonrigorous fashion in the introduction.

Example. Let $\mathcal{A}=\{n(n+2): n \leq x\}$ and let $\mathcal{P}=\{p \leq z\}$. As pointed out in (6.10), for this choice of $A$ and $\mathcal{P}, S(\mathcal{A}, \mathcal{P})$ is related to the twin-prime counting function $\pi_{2}(x)$. In order to decide on a reasonable choice of $X$ and $\alpha$ in this situation, let us attempt to get a feel for the numbers $A_{d}$. The condition that $d$ divides $n(n+2)$ is a condition on $n$ modulo $d$, so we set

$$
\begin{equation*}
\nu(d):=\#\{n(\bmod d): n(n+2) \equiv 0(\bmod d)\} . \tag{6.18}
\end{equation*}
$$

Then each block of $d$ consecutive integers contains precisely $\nu(d)$ solutions of the congruence $n(n+2) \equiv 0(\bmod d)$. Hence $A_{d} \approx(x / d) \nu(d)$, which suggests that we choose $X=x$ and $\alpha(d)=\nu(d) / d$. (Note that $\nu$, and hence $\alpha$, is multiplicative by the Chinese remainder theorem.) In fact, since the interval $[1, x]$ contains the first $\lfloor x / d\rfloor$ blocks of $d$ consecutive natural numbers, and is contained in the first $\lceil x / d\rceil$ such blocks, with this choice of $X$ and $\alpha$ we have

$$
\lfloor x / d\rfloor \nu(d) \leq A_{d} \leq\lceil x / d\rceil \nu(d), \quad \text { so that } \quad|r(d)|=\left|A_{d}-x \nu(d) / d\right| \leq \nu(d) .
$$

We now apply Theorem 6.2, with $z$ a function of $x$ tending slowly to infinity. The coefficient of $X=x$ in the main term of Theorem 6.2 is

$$
\prod_{p \leq z}(1-\alpha(p))=\frac{1}{2} \prod_{2<p \leq z}\left(1-\frac{2}{p}\right)=\left(2 \prod_{2<p \leq z} \frac{1-\frac{2}{p}}{\left(1-\frac{1}{p}\right)^{2}}\right) \prod_{p \leq z}\left(1-\frac{1}{p}\right)^{2}
$$

Estimating the last product here by Mertens' theorem, we find that the main term is asymptotic to

$$
2 C_{2} e^{-2 \gamma} x /(\log z)^{2} .
$$

The error term is bounded by

$$
\sum_{d \mid P} \nu(d)=\prod_{p \leq z}(1+\nu(p)) \leq 3^{\pi(z)} \leq 3^{z}
$$

which is negligible in comparison with the main term if (e.g.) $z=\frac{1}{2} \log x$. We will return to this example in $\S 4.4$.

Example. Here is an example different from those alluded to in the introduction, due to Nagell [ $\mathbf{N a g} 22$ ]. Let $\pi_{T^{2}+1}(x)$ denote the number of $n \leq x$ for which $n^{2}+1$ is prime. Let $\mathcal{A}=\left\{n^{2}+1: n \leq x\right\}$ and let $\mathcal{P}$ be the set of all primes. Then for any choice of positive numbers $x$ and $z$, we have

$$
\begin{equation*}
\pi_{T^{2}+1}(x) \leq S(\mathcal{A}, \mathcal{P}, z)+z^{1 / 2} \tag{6.19}
\end{equation*}
$$

The congruence $n^{2}+1 \equiv 0(\bmod d)$ is satisfied precisely when $n$ falls into one of $\nu(d)$ (say) residue classes modulo $d$. As in the preceding example, this suggests we take $X=x$ and $\alpha(d)=\nu(d) / d$; with this choice of $X$ and $\alpha$, the numbers $r(d)$ defined by (6.16) satisfy $|r(d)| \leq \nu(d)$.

Now $\nu(2)=1$, while if $p$ is an odd prime, $\nu(p)=0$ or 2 , depending on whether $p \equiv 3(\bmod 4)$ or $1(\bmod 4)$, respectively. So by $(6.19)$ and Theorem 6.2, if $x>0$ and $z \geq 2$, then

$$
\begin{aligned}
\pi_{T^{2}+1}(x) & \leq S(\mathcal{A}, \mathcal{P}, z)+z^{1 / 2} \\
& \leq \frac{1}{2} x \prod_{\substack{p \leq z \\
p \equiv 1(\bmod 4)}}\left(1-\frac{2}{p}\right)+O\left(\sum_{d \mid P_{z}} \nu(d)\right)+z^{1 / 2} .
\end{aligned}
$$

To understand the main term, note that

$$
\prod_{\substack{p \leq z \\ p \equiv 1(\bmod 4)}}\left(1-\frac{2}{p}\right) \leq \exp \left(-2 \sum_{\substack{p \leq z \\ p \equiv 1(\bmod 4)}} \frac{1}{p}\right) \ll \frac{1}{\log z} .
$$

Here we have used that $\sum_{p \leq z, p \equiv 1(\bmod 4)} \frac{1}{p}=\frac{1}{2} \log \log z+O(1)$, which follows by partial summation from the results of Chapter 4 . For the $O$-term, we have $\sum_{d \mid P_{z}} \nu(d)=\prod_{p \leq z}(1+\nu(p))<3^{z}$. Inserting these estimates above and choosing $z=\frac{1}{2} \log x$, we find that $\pi_{T^{2}+1}(x) \ll x / \log \log x$. In particular, the set of numbers $n$ for which $n^{2}+1$ is prime has density zero.

The following simple consequence of Theorem 6.2 is often useful:

Corollary 6.3. Let $\mathcal{P}$ be a set of prime numbers, and let $M(\mathcal{P})$ denote the set of $n \in \mathbf{N}$ divisible by some prime $p \in \mathcal{P}$. Then $M(\mathcal{P})$ has asymptotic density $1-\prod_{p \in \mathcal{P}}(1-1 / p)$. In particular, $M(\mathcal{P})$ has density 1 precisely when $\sum_{p \in \mathcal{P}} p^{-1}$ diverges.

Proof. Let $M^{\prime}=\mathbf{N} \backslash M(\mathcal{P})$ be the set of natural numbers $n$ divisible by none of the elements of $\mathcal{P}$, and write $M^{\prime}(x)$ for the associated counting function. Put $\mathcal{A}:=\{n \leq x\}$. Then for any choice of $z$, we have

$$
\begin{equation*}
M^{\prime}(x) \leq S(\mathcal{A}, \mathcal{P}, z) \tag{6.20}
\end{equation*}
$$

The right-hand side will be estimated with the aid of Theorem 6.2. We take $X=x$ and let $\alpha$ be the multiplicative function with $\alpha(n):=1 / n$ for every $n \in \mathbf{N}$. With this choice of $X$ and $\alpha$, we have $|r(d)| \leq 1$ for every $d \mid P$. Now put $z=\log x$. By Theorem 6.2,

$$
\begin{align*}
S(\mathcal{A}, \mathcal{P}, z) & =x \prod_{\substack{p \in \mathcal{P} \\
p \leq \log x}}(1-1 / p)+O\left(2^{\log z}\right) \\
& =(C+o(1)) x, \quad \text { where } \quad C:=\prod_{p \in \mathcal{P}}(1-1 / p) . \tag{6.21}
\end{align*}
$$

If $C=0$, then we obtain from (6.20) and (6.21) that $M^{\prime}$ has density zero, so that $M(\mathcal{P})$ has density 1 , which is the assertion of the corollary in this case. If $C \neq 0$, then $\sum_{p \in \mathcal{P}} p^{-1}$ converges, and so

$$
\begin{aligned}
M^{\prime}(x) & \geq S(\mathcal{A}, \mathcal{P}, z)-\sum_{\substack{p \in \mathcal{P} \\
p>z}} \frac{x}{p} \\
& =(C+o(1)) x+o(x)=(C+o(1)) x .
\end{aligned}
$$

With (6.21), this shows that $M^{\prime}$ has asymptotic density $C$, so that $M(\mathcal{P})$ has density $1-C$, as desired.

Suppose that in Corollary 6.3 we take $\mathcal{P}$ to be the entire set of prime numbers. Then $M(\mathcal{P})$ consists of every natural number $n>1$ and so has density 1. Thus $\prod_{p \in \mathcal{P}}(1-1 / p)=0$. This gives another proof of Euler's result from Chapter 1 that $\sum_{p} \frac{1}{p}$ diverges (cf. [Pin09]).
3.3. Three applications. We pause to give three further applications of Corollary 6.3. None of the results we prove are the best of their kind, but the proofs are simple and the statements fairly striking.

Theorem 6.4. Each of the following sets has density zero:
(i) the set of integers $n>1$ for which the equation

$$
\begin{equation*}
4 / n=1 / a+1 / b+1 / c \tag{6.22}
\end{equation*}
$$

has no solution in positive integers $a, b, c$,
(ii) the set of natural numbers expressible as a sum of two squares, (iii) the set of odd perfect numbers.

Remark. The set in (iii) is famously conjectured to be empty; we discuss this conjecture at length in Chapter 8. Erdős \& Straus (see [Erd50a]) believe that the same holds for the set in (i), that is, that $4 / n$ can always be written as a sum of three unit fractions (for $n>1$ ). For example,

$$
\frac{4}{301}=\frac{1}{76}+\frac{1}{7626}+\frac{1}{87226188}
$$

Of course, the analogous conjecture is trivial if "three" is replaced by "four". It has been verified by computer that the set in (i) contains no $n \leq 10^{14}$.

As regards (ii), Landau [Lan08] has proved that the number of $n \leq x$ expressible as a sum of two squares is

$$
\sim \frac{1}{\sqrt{2}}\left(\prod_{p \equiv 3(\bmod 4)}\left(1-\frac{1}{p^{2}}\right)\right)^{-1 / 2} \frac{x}{\sqrt{\log x}}
$$

The simplest proof of Landau's result seems to be that of Selberg [Sel91, pp. 183-185].

Lemma 6.5. The set of positive integers divisible by no prime $p \equiv 3(\bmod 4)$ has density 0 .

Proof. From Chapter 4, we have that $\sum_{p \leq x}^{\prime} p^{-1} \log p=\frac{1}{2} \log x+O(1)$, where the ${ }^{\prime}$ indicates that the sum is restricted to primes $p \equiv 3(\bmod 4)$. So by partial summation, $\sum_{p \leq x}^{\prime} p^{-1} \sim \frac{1}{2} \log \log x$. In particular, $\sum_{p}^{\prime} p^{-1}$ diverges. So the result follows from Corollary 6.3.

Proof of Theorem 6.4(i). It suffices to show that (6.22) is solvable if $n$ possesses a prime divisor $p=4 k-1 \equiv 3(\bmod 4)$. In this case write $n=(4 k-1) q$. Then

$$
\frac{4}{n}=\frac{4}{q(4 k-1)}=\frac{1}{2 q k}+\frac{1}{2 q k}+\frac{1}{q\left(4 k^{2}-k\right)} .
$$

This argument also shows that $4 / n$ can almost always be written as a sum of two unit fractions, since $1 /(2 q k)+1 /(2 q k)=1 /(q k)$.

Proof of Theorem 6.4(ii). Let $R(x)$ be the number of $n \leq x$ which can be written as a sum of two squares, and let $A(x)$ be the number of $n \leq x$ which have a primitive representation of this form, i.e., a representation as a sum of two coprime squares. As shown by Euler, the $n$ counted by $A(x)$
are precisely those divisible by neither 4 nor any prime $p \equiv 3(\bmod 4)$. Moreover,

$$
\begin{equation*}
R(x) \leq A\left(\frac{x}{1^{2}}\right)+A\left(\frac{x}{2^{2}}\right)+A\left(\frac{x}{3^{2}}\right)+\cdots . \tag{6.23}
\end{equation*}
$$

By Lemma 6.5, we have $A(x)=o(x)$. Now given $\epsilon>0$, choose an $N \in \mathbf{N}$ for which $A(x)<\epsilon x / 4$ whenever $x>N$. Thinking of $x$ as large, we split the sum in (6.23) into two parts according to whether $x / k^{2}>N$ or $x / k^{2} \leq N$. The first of the two resulting sums is bounded by

$$
\frac{1}{4} \epsilon \sum_{k \geq 1} \frac{x}{k^{2}}=\epsilon \frac{\zeta(2)}{4} x<\frac{\epsilon}{2} x .
$$

Every term in the second sum is bounded by $A(N)$, and there are no more than $\sqrt{x}$ nonzero terms. Thus,

$$
R(x) \leq \epsilon x / 2+A(N) \sqrt{x}<\epsilon x
$$

for large $x$. As $\epsilon>0$ was arbitrary, it follows that $R(x)=o(x)$.
Proof of Theorem 6.4(iii). It has been known since Euler that every odd perfect number $n$ can be written in the form $p a^{2}$, where $p \equiv 1(\bmod 4)$ is prime. (We will prove a stronger version of this result in Chapter 8; see Theorem 8.2 there.) Since such integers are sums of two squares, the result follows from that of part (ii).

## 4. Brun's pure sieve

In the derivation of Legendre's formula for $\pi(x, z)$ given in $\S 1.2$ above, we begin with the total number of positive integers not exceeding $x$. For each prime $p \leq z$, we take away the number of multiples of $p$. Then, for each pair of primes $p<q \leq z$, we add back the number $n$ divisible by both $p$ and $q$. Continuing we eventually converge on the exact value of $\pi(x, z)$. It is intuitively clear (and we will prove it below) that after each even (addition) step what we have is an overestimate for $\pi(x, z)$, and after each odd (subtraction) step we have an underestimate. A suitable generalization of this fact forms the heart of Brun's pure sieve.
4.1. Preparation. To prove the appropriate generalization, it is convenient to first establish a technical lemma on alternating sums of symmetric functions.

If $a_{1}, \ldots, a_{n}$ is a (possibly empty) sequence of $n \geq 0$ elements belonging to a commutative ring, we define (for $k \geq 0$ ) the $k$ th elementary symmetric function $\sigma_{k}\left(a_{1}, \ldots, a_{n}\right)$ as the sum of all possible $\binom{n}{k}$ products of the $a_{i}$ taken $k$ at a time. We adopt the usual conventions about empty sums and
products, so that for $n=0$, we have $\sigma_{0}=1$ and $\sigma_{k}=0$ for $k>0$. To take a less pathological example, when $n=2$, one has

$$
\sigma_{0}\left(a_{1}, a_{2}\right)=1, \quad \sigma_{1}\left(a_{1}, a_{2}\right)=a_{1}+a_{2}, \quad \sigma_{2}\left(a_{1}, a_{2}\right)=a_{1} a_{2},
$$

and $\sigma_{k}\left(a_{1}, a_{2}\right)=0$ for $k>2$. The following lemma can be found, e.g., in [Hoo94]:
Lemma 6.6. Suppose $0 \leq a_{1}, \ldots, a_{n} \leq 1$, where $n$ is nonnegative. Then

$$
\begin{equation*}
\sum_{k=0}^{m}(-1)^{k} \sigma_{k}\left(a_{1}, \ldots, a_{n}\right)-\prod_{j=1}^{n}\left(1-a_{j}\right) \tag{6.24}
\end{equation*}
$$

is nonnegative or nonpositive according to whether $m$ is even or odd, respectively.
Remark. Note that (6.24) vanishes when $m \geq n$.
Proof. We induct on the length $n$ of the sequence. When $n=0$, the product $P:=\prod_{i=1}^{n}\left(1-a_{i}\right)$ appearing in (6.24) is empty, so equal to 1 , while

$$
\sum_{k=0}^{m}(-1)^{k} \sigma_{k}=1-0+0-\cdots \pm 0=1
$$

Hence (6.24) vanishes for every $m$, confirming the result in this case. Now assume that the result holds for each sequence of $n$ real numbers in $[0,1]$ and each $m$, and consider an arbitrary sequence $0 \leq a_{1}, \ldots, a_{n+1} \leq 1$ of length $n+1$. By the induction hypothesis, it suffices to prove that

$$
\begin{align*}
& \left(\sum_{k=0}^{m}(-1)^{k} \sigma_{k}\left(a_{1}, \ldots, a_{n+1}\right)-\prod_{i=1}^{n+1}\left(1-a_{i}\right)\right)  \tag{6.25}\\
& -\left(\sum_{k=0}^{m}(-1)^{k} \sigma_{k}\left(a_{1}, \ldots, a_{n}\right)-\prod_{i=1}^{n}\left(1-a_{i}\right)\right)
\end{align*}
$$

is nonnegative or nonpositive according to whether $m$ is even or odd respectively. This is easily seen to hold for $m=0$, since then (6.25) simplifies to $P a_{n+1}$, which is nonnegative. When $m>0$, we can rewrite (6.25) as

$$
\begin{aligned}
\sum_{k=1}^{m}(-1)^{k}\left(\sigma_{k}\left(a_{1}, \ldots, a_{n+1}\right)-\right. & \left.\sigma_{k}\left(a_{1}, \ldots, a_{n}\right)\right)+P a_{n+1} \\
& =\sum_{k=1}^{m}(-1)^{k} a_{n+1} \sigma_{k-1}\left(a_{1}, \ldots, a_{n}\right)+P a_{n+1} \\
& =a_{n+1}\left(P-\sum_{k=0}^{m-1}(-1)^{k} \sigma_{k}\left(a_{1}, \ldots, a_{n}\right)\right)
\end{aligned}
$$

The claim in this case now follows from the induction hypothesis.

An important special case occurs when $n \in \mathbf{N}$ and $a_{1}=a_{2}=\cdots=a_{n}=$

1. Then $\prod_{i=1}^{n}\left(1-a_{i}\right)=(1-1)^{n}=0$, while $\sigma_{k}(1, \ldots, 1)=\binom{n}{k}$. So from Lemma 6.6 we obtain the following:
Lemma 6.7. Let $n$ be a positive integer. Then the alternating sum

$$
\sum_{k=0}^{m}(-1)^{k}\binom{n}{k}
$$

is nonnegative or nonpositive according to whether $m$ is even or odd.
Remark. While Lemma 6.6 will be important in our treatment of the BrunHooley sieve, for Brun's pure sieve we only need Lemma 6.7. Thus it is of interest that Lemma 6.7 admits a simple proof independent of Lemma 6.6: Indeed, by induction on $m$, one easily finds that

$$
\begin{equation*}
\sum_{k=0}^{m}(-1)^{k}\binom{n}{k}=(-1)^{m}\binom{n-1}{m} \tag{6.26}
\end{equation*}
$$

which makes Lemma 6.7 obvious. Alternatively, (6.26) follows by comparing the coefficient of $x^{m}$ in both sides of the power series identity $(1-x)^{n-1}=$ $(1-x)^{-1}(1-x)^{n}$.

Lemma 6.7 implies the following variant of Theorem 6.1:
Theorem 6.8 (Bonferroni inequalities). Let $X$ be a nonempty, finite set of $N$ objects, and let $P_{1}, \ldots, P_{r}$ be properties that elements of $X$ may have. For each subset $I \subset\{1,2, \ldots, r\}$, let $N(I)$ denote the number of elements of $X$ that have each of the properties indexed by the elements of $I$. Let $N_{0}$ denote the number of elements of $X$ with none of these properties. Then if $m$ is a nonnegative even integer,

$$
\begin{equation*}
N_{0} \leq \sum_{k=0}^{m}(-1)^{k} \sum_{\substack{I \subset\{1,2, \ldots, r\} \\|I|=k}} N(I) \tag{6.27}
\end{equation*}
$$

while if $m$ is a nonnegative odd integer,

$$
\begin{equation*}
N_{0} \geq \sum_{k=0}^{m}(-1)^{k} \sum_{\substack{I \subset\{1,2, \ldots, r\} \\|I|=k}} N(I) \tag{6.28}
\end{equation*}
$$

Proof. Suppose that $x \in X$ has exactly $l$ of the properties $P_{1}, \ldots, P_{r}$. If $l=0$, then $x$ is counted once by both $N_{0}$ and the common right-hand side of (6.27) and (6.28) (corresponding to $I=\emptyset$ ). If $l \geq 1$, then $x$ is not counted at all by $N_{0}$, and is counted by this right-hand sum with weight

$$
\sum_{k=0}^{m}(-1)^{k}\binom{l}{k} \begin{cases}\geq 0 & \text { if } m \text { is even } \\ \leq 0 & \text { otherwise }\end{cases}
$$

Summing over $x \in X$ gives the theorem.

### 4.2. A working version.

Corollary 6.9 (Brun's pure sieve, general form). With the notation of §2, we have for every nonnegative even integer $m$,

$$
\sum_{d \mid P, \omega(d) \leq m-1} \mu(d) A_{d} \leq S(\mathcal{A}, \mathcal{P}) \leq \sum_{d \mid P, \omega(d) \leq m} \mu(d) A_{d} .
$$

Proof. As in the proof of Theorem 6.2, let $p_{1}, \ldots, p_{r}$ be a list of the primes $p \in \mathcal{P}$, and let $P_{i}$ be the property of being divisible by $p_{i}$. We aim to estimate the number $S(\mathcal{A}, \mathcal{P})$ of elements of $\mathcal{A}$ possessing none of the $P_{i}$. The upper bound for $S(\mathcal{A}, \mathcal{P})$ in the corollary is just (6.27). If $m=0$, then the lower bound is trivial, while if $m>0$, then $m-1$ is a nonnegative odd integer, and the lower bound follows from (6.28).

To obtain a result suitable for applications, we substitute $A_{d}=X \alpha(d)+$ $r(d)$. With a bit of manipulation, we arrive at the following theorem:

Theorem 6.10 (Brun's pure sieve). For every even integer $m \geq 0$,

$$
S(\mathcal{A}, \mathcal{P})=X \prod_{p \in \mathcal{P}}(1-\alpha(p))+O\left(\sum_{d \mid P, \omega(d) \leq m}|r(d)|\right)+O\left(X \sum_{d \mid P, \omega(d) \geq m} \alpha(d)\right) .
$$

Here the implied constants are absolute.

Proof. From Corollary 6.9,

$$
\begin{aligned}
S(\mathcal{A}, \mathcal{P}) & =\sum_{\substack{d \mid P \\
\omega(d) \leq m}} \mu(d) A_{d}+O\left(\sum_{\substack{d \mid P \\
\omega(d)=m}} A_{d}\right) \\
& =\sum_{\substack{d \mid P \\
\omega(d) \leq m}} \mu(d)(X \alpha(d)+r(d))+O\left(\sum_{\substack{d \mid P \\
\omega(d)=m}} A_{d}\right) \\
& =X \sum_{\substack{d \mid P \\
\omega(d) \leq m}} \mu(d) \alpha(d)+O\left(\sum_{\substack{d \mid P \\
\omega(d) \leq m}}|r(d)|\right)+O\left(\sum_{\substack{d \mid P \\
\omega(d)=m}} A_{d}\right) .
\end{aligned}
$$

Writing $A_{d}=X \alpha(d)+r(d)$, we see that the last of these error terms is

$$
\ll X \sum_{d \mid P, \omega(d)=m} \alpha(d)+\sum_{d \mid P, \omega(d)=m}|r(d)| ;
$$

hence,

$$
\begin{align*}
& S(\mathcal{A}, \mathcal{P})  \tag{6.29}\\
& \quad=X \sum_{\substack{d \mid P \\
\omega(d) \leq m}} \mu(d) \alpha(d)+O\left(\sum_{\substack{d \mid P \\
\omega(d) \leq m}}|r(d)|\right)+O\left(X \sum_{\substack{d \mid P \\
\omega(d)=m}} \alpha(d)\right) .
\end{align*}
$$

In order to factor the sum appearing in the main term, we extend the sum to all $d \mid P$; the main term can then be expressed as $X \prod_{p \in \mathcal{P}}(1-\alpha(p))$, but we have introduced a new error of

$$
\ll X \sum_{d \mid P, \omega(d)>m} \alpha(d) .
$$

If this is combined with the last error term of (6.29), we find that

$$
S(\mathcal{A}, \mathcal{P})=X \prod_{p \in \mathcal{P}}(1-\alpha(p))+O\left(\sum_{d \mid P, \omega(d) \leq m}|r(d)|\right)+O\left(X \sum_{d \mid P, \omega(d) \geq m} \alpha(d)\right)
$$

exactly as the theorem asserts.
4.3. Application to the twin prime problem. The most famous application of Brun's pure sieve is Brun's own 1919 contribution [Bru19a] to the twin prime problem:

Theorem 6.11. As $x \rightarrow \infty$,

$$
\pi_{2}(x) \ll \frac{x}{(\log x)^{2}}(\log \log x)^{2}
$$

The upper bound differs from what is expected by a factor of $(\log \log x)^{2}$. We shall later remedy this defect. Nevertheless, it is worth noting that the estimate of Theorem 6.11 is already sharp enough to imply the following striking result:

Corollary 6.12. If there are infinitely many primes $p$ such that $p+2$ is also prime, then the sum

$$
\sum_{p} \frac{1}{p}
$$

taken over all such primes, converges.
Proof. By Theorem 6.13, $\pi_{2}(x) \ll x /(\log x)^{3 / 2}$ as $x \rightarrow \infty$. It follows that the same estimate holds, with perhaps a different implied constant, in the range $x \geq 3$. Letting $p_{n}$ denote the $n$th prime $p$ for which $p+2$ is also prime, we see that for $n \geq 1$,

$$
n=\pi_{2}\left(p_{n}\right) \ll p_{n} /\left(\log p_{n}\right)^{3 / 2}
$$

so that

$$
p_{n} \gg n\left(\log p_{n}\right)^{3 / 2} \geq \frac{1}{2}(n+1)(\log (n+1))^{3 / 2} .
$$

The comparison and integral tests together now imply that $\sum_{n=1}^{\infty} p_{n}^{-1}$ converges, which is the assertion of the corollary.

Remark. For historical reasons, in place of the series appearing in Corollary 6.12 one usually sees the slight variant

$$
\left(\frac{1}{3}+\frac{1}{5}\right)+\left(\frac{1}{5}+\frac{1}{7}\right)+\left(\frac{1}{11}+\frac{1}{13}\right)+\cdots .
$$

Of course this series converges (by comparison with that of the corollary), and its value $B$ is known as Brun's constant. Computing the value of $B$ to any precision seems to be difficult; while constants like $\pi$ and $e$ are known to billions of decimal digits, the sharpest known bounds on $B$ are (roughly)

$$
1.830<B<2.347 .
$$

Thus we do not know $B$ to even one significant digit! The lower bound here is due to Sebah $[\mathbf{S G}]$, who computed all the twin prime pairs up to $10^{16}$ and summed their reciprocals. The upper bound is due to Crandall \& Pomerance ([CP05, pp. 16-17], see also [Kly07, Chapter 3]), who bound the sum of the twin prime pairs past $10^{16}$ using an explicit upper estimate of Riesel and Vaughan [RV83] for the number of twin prime pairs. Much sharper estimates for Brun's constant are available if one assumes a suitable quantitative version of the twin prime conjecture; e.g., it is plausible that

$$
B=1.902160583121 \pm 4.08 \times 10^{-8}
$$

This last estimate is taken from the Ph.D. thesis of Klyve [Kly07], which the reader should consult for references to earlier work.

With $\mathcal{A}:=\{n(n+2): n \leq x\}$ and $\mathcal{P}:=\{p \leq z\}$, put $\pi_{2}(x, z):=S(\mathcal{A}, \mathcal{P})$. Theorem 6.11 is an easy consequence of the following estimate:

Theorem 6.13. Suppose $z=z(x) \rightarrow \infty$ as $x \rightarrow \infty$ and that $z(x) \leq$ $x^{1 /(20 \log \log x)}$ for all large $x$. Then $\pi_{2}(x, z) \sim 2 C_{2} e^{-2 \gamma} x /(\log z)^{2}$ as $x \rightarrow \infty$, where $C_{2}$ is the twin prime constant.

Proof of Theorem 6.11 assuming Theorem 6.13. Relation (6.10) tells us that $\pi_{2}(x) \leq z+\pi_{2}(x, z)$. Take $z=x^{1 /(20 \log \log x)}$. Theorem 6.13 implies that as $x \rightarrow \infty$,

$$
\pi_{2}(x) \ll x^{1 /(20 \log \log x)}+\frac{x}{(\log x)^{2}}(\log \log x)^{2} \ll \frac{x}{(\log x)^{2}}(\log \log x)^{2} .
$$

4.4. Proof of Theorem 6.13. Estimates for $\pi_{2}(x, z)$ were discussed in the second example of $\S 3.2$; the difference here is that we now have Brun's pure sieve at our disposal. As in that example, we take $X=x$ and $\alpha(d)=\nu(d) / d$, where $\nu$ is defined by (6.18). Then $|r(d)| \leq \nu(d) \leq 2^{\omega(d)}$ for all $d$. So by Theorem 6.10,

$$
\begin{align*}
& \pi_{2}(x, z)=  \tag{6.30}\\
& \quad x \prod_{p \leq z}(1-\alpha(p))+O\left(\sum_{d \mid P, \omega(d) \leq m} 2^{\omega(d)}\right)+O\left(x \sum_{d \mid P, \omega(d) \geq m} \alpha(d)\right),
\end{align*}
$$

for each even number $m \geq 0$. We take

$$
m:=10\lfloor\log \log z\rfloor .
$$

Note that as $x$ goes to infinity, so does $z$ and hence also $m$. In $\S 3.2$, we calculated that the main term of (6.30) is asymptotic to

$$
2 C_{2} e^{-2 \gamma} x /(\log z)^{2}
$$

as $x \rightarrow \infty$. So to prove Theorem 6.13, it is enough to establish the following two estimates:
(i) With $E_{1}:=\sum_{d \mid P, \omega(d) \leq m} 2^{\omega(d)}$, we have $E_{1}=o\left(x /(\log z)^{2}\right)$.
(ii) With $E_{2}:=x \sum_{d \mid P, \omega(d) \geq m} \alpha(d)$, we have $E_{2}=o\left(x /(\log z)^{2}\right)$.

Proof of (i). For large $x$,

$$
\begin{aligned}
E_{1}=\sum_{d \mid P, \omega(d) \leq m} 2^{\omega(d)} & =\sum_{k=0}^{m} 2^{k}\binom{\pi(z)}{k} \leq \sum_{k=0}^{m}(2 \pi(z))^{k} \\
& \leq \sum_{k=-\infty}^{m}(2 \pi(z))^{k}=(2 \pi(z))^{m} \frac{1}{1-\frac{1}{2 \pi(z)}} \\
& \leq 2(2 \pi(z))^{m} \leq 2 z^{m},
\end{aligned}
$$

since $\pi(z) \leq z / 2$ for large $x$. Hence

$$
E_{1} \leq 2 z^{10 \log \log z} \leq 2 z^{10 \log \log x} \leq 2 x^{1 / 2}
$$

This upper bound is certainly $o\left(x /(\log z)^{2}\right)$, since as $x \rightarrow \infty$,

$$
\frac{x^{1 / 2}}{x /(\log z)^{2}} \leq \frac{x^{1 / 2}}{x /(\log x)^{2}}=\frac{(\log x)^{2}}{x^{1 / 2}} \rightarrow 0 .
$$

Proof of (ii). We can write $E_{2}=x \sum_{k \geq m} \sum_{d \mid P, \omega(d)=k} \alpha(d)$. For the inner sum we have

$$
\sum_{\substack{d \mid P \\ \omega(d)=k}} \alpha(d)=\sum_{p_{1}<p_{2}<\cdots<p_{k} \leq z} \alpha\left(p_{1}\right) \alpha\left(p_{2}\right) \cdots \alpha\left(p_{k}\right) \leq \frac{1}{k!}\left(\sum_{p \leq z} \alpha(p)\right)^{k}
$$

Here the upper bound comes from the multinomial theorem: In the expansion of $\left(\sum_{p \leq z} \alpha(p)\right)^{k}$, every term $\alpha\left(p_{1}\right) \cdots \alpha\left(p_{k}\right)$ appears with coefficient $k!$. From Mertens' first theorem, we have $\sum_{p \leq z} p^{-1} \leq \log \log z+c$ for $z \geq 3$, where $c$ is an absolute constant. Since $\alpha(p) \leq 2 / p$ for every prime $p$,

$$
\begin{equation*}
\sum_{k \geq m} \frac{1}{k!}\left(\sum_{p \leq z} \alpha(p)\right)^{k} \leq \sum_{k \geq m} \frac{1}{k!}(2 \log \log z+2 c)^{k} . \tag{6.31}
\end{equation*}
$$

The ratio of the $(k+1)$ th term in the right-hand series to the $k$ th is given by

$$
\frac{2 \log \log z+2 c}{k+1} \leq \frac{2 \log \log z+2 c}{10\lfloor\log \log z\rfloor+1} \leq 1 / 2,
$$

for large enough $z$, and hence also for large enough $x$. So, for such $x$ the right-hand sum in (6.31) is bounded by twice its first term. Because

$$
e^{m}=1+m+m^{2} / 2!+m^{3} / 3!+\cdots \geq m^{m} / m!,
$$

we have $m!\geq(m / e)^{m}$, so that

$$
\sum_{k \geq m} \frac{1}{k!}(2 \log \log z+2 c)^{k} \leq 2\left(\frac{2 e \log \log z+2 c e}{m}\right)^{m}
$$

Since $m=10\lfloor\log \log z\rfloor$, the parenthetical expression on the right is eventually smaller than any constant exceeding $2 e / 10$; in particular, it is eventually smaller than $3 / 5$. It follows that for large $x$,

$$
\begin{aligned}
E_{2} \leq 2 x(3 / 5)^{m} & =2 x(3 / 5)^{10\lfloor\log \log z\rfloor} \\
& \ll x(3 / 5)^{10 \log \log z} \ll x /(\log z)^{5},
\end{aligned}
$$

since $10 \log \frac{3}{5}<-5$. So $E_{2}=o\left(x /(\log z)^{2}\right)$.

## 5. The Brun-Hooley sieve

5.1. The sifting function perspective. Before we discuss the BrunHooley method, it is worthwhile for us to revisit some of the earlier results of this chapter from a a slightly different perspective. Keeping the notation of $\S 2$, we introduce the sifting function

$$
s(n):= \begin{cases}1 & \text { if } \operatorname{gcd}(n, P)=1  \tag{6.32}\\ 0 & \text { otherwise }\end{cases}
$$

Then

$$
\begin{equation*}
S(\mathcal{A}, \mathcal{P})=\sum_{a \in A} s(a) \tag{6.33}
\end{equation*}
$$

Since $\sum_{d \mid m} \mu(d)$ vanishes for each natural number $m>1$, the sifting function $s(n)$ has the following important representation:

$$
\begin{equation*}
s(n)=\sum_{d|n, d| P} \mu(d) . \tag{6.34}
\end{equation*}
$$

Substituting this into (6.33) and interchanging the order of summation, we easily arrive at Theorem 6.2 (the sieve of Eratosthenes-Legendre). In the same way, Brun's pure sieve is a consequence of the following lemma:
Lemma 6.14. Let $n$ be a natural number. The expression

$$
\begin{equation*}
\sum_{\substack{d|n, d| P \\ \omega(d) \leq m}} \mu(d)-\sum_{d|n, d| P} \mu(d) \tag{6.35}
\end{equation*}
$$

is nonnegative or nonpositive according to whether the integer $m \geq 0$ is even or odd.

The proof of Lemma 6.14 is essentially the one already given for the Bonferroni inequalities. Namely, if we suppose that $n$ is divisible by exactly $l$ primes $p \in \mathcal{P}$, then by Lemma 6.6,

$$
\sum_{\substack{d|n, d| P \\ \omega(d) \leq m}} \mu(d)=\sum_{k=0}^{m}(-1)^{k}\binom{l}{k} \begin{cases}=1 & \text { if } l=0 \text { (i.e., if } \operatorname{gcd}(n, P)=1) \\ \geq 0 & \text { if } l \geq 1, m \text { even, } \\ \leq 0 & \text { if } l \geq 1, m \text { odd. }\end{cases}
$$

For later use we note the following consequence of Lemma 6.14:
Lemma 6.15. If $n$ is a natural number and $m \geq 0$ is even, then

$$
0 \leq \sum_{\substack{d|n, d| P \\ \omega(d) \leq m}} \mu(d)-\sum_{d|n, d| P} \mu(d) \leq \sum_{\substack{d|n, d| P \\ \omega(d)=m+1}} 1 .
$$

5.2. The upper bound. The Brun-Hooley method takes two forms, depending on whether we are after upper or lower bounds. Here we describe the simpler upper bound method. We suppose the sifting set $\mathcal{P}$ to be partitioned into $r$ disjoint sets, say $\mathcal{P}=\dot{\bigcup}_{j=1}^{r} \mathcal{P}_{j}$. Then $n$ is divisible by no prime $p \in \mathcal{P}$ precisely when $n$ is divisible by no prime $p \in \mathcal{P}_{j}$ for every $1 \leq j \leq r$. Consequently, setting $P_{j}:=\prod_{p \in \mathcal{P}_{j}} p$, and invoking Lemma 6.14 (with $\mathcal{P}_{j}, P_{j}$ in place of $\mathcal{P}, P$ ) we see that

$$
\begin{aligned}
s(n)=\sum_{d|n, d| P} \mu(d) & =\prod_{j=1}^{r} \sum_{d_{j}\left|n, d_{j}\right| P_{j}} \mu\left(d_{j}\right) \\
& \leq \prod_{\substack{j=1 \\
r}} \sum_{\substack{d_{j}\left|n, d_{j}\right| P_{j} \\
\omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{j}\right),
\end{aligned}
$$

for any choice of nonnegative even integers $m_{1}, \ldots, m_{r}$. Referring to (6.33), we obtain the upper bound

$$
\begin{aligned}
& S(\mathcal{A}, \mathcal{P}) \leq \sum_{\substack{d_{1}, \ldots, d_{r} \leq m_{j} \\
d_{j} \mid P_{j}, \omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{1}\right) \cdots \mu\left(d_{r}\right) A_{d_{1} \cdots d_{r}} \\
& =X \sum_{\substack{d_{1}, \ldots, d_{r} \\
d_{j} \mid P_{j}, \omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{1}\right) \cdots \mu\left(d_{r}\right) \alpha\left(d_{1}\right) \cdots \alpha\left(d_{r}\right) \\
&
\end{aligned}
$$

Hence $S(\mathcal{A}, \mathcal{P})$ is bounded above by

$$
\begin{equation*}
X \prod_{j=1}^{r} \sum_{\substack{d_{j} \mid P_{j} \\ \omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{j}\right) \alpha\left(d_{j}\right)+\sum_{\substack{d_{1}, \ldots, d_{r} \\ d_{j} \mid P_{j}, \omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{1}\right) \cdots \mu\left(d_{r}\right) r\left(d_{1} \cdots d_{r}\right) . \tag{6.37}
\end{equation*}
$$

This is the upper bound of the Brun-Hooley method. To facilitate applications, we replace the first term of (6.37), which we think of as the main term, with something more easily compared with $X \prod_{p \in \mathcal{P}}(1-\alpha(p))$. This can be accomplished by replacing the $j$ th term of the product in (6.37) with something more easily compared with $\prod_{p \in \mathcal{P}_{j}}(1-\alpha(p))$. For this, we utilize Lemma 6.6, which implies that for each $1 \leq j \leq r$,

$$
0 \leq \sum_{\substack{d_{j} \mid P_{j} \\ \omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{j}\right) \alpha\left(d_{j}\right)-\prod_{p \in \mathcal{P}_{j}}(1-\alpha(p)) \leq \sum_{\substack{d_{j} \mid P_{j} \\ \omega\left(d_{j}\right)=m_{j}+1}} \alpha\left(d_{j}\right)
$$

Thus, if we set

$$
\begin{equation*}
\prod^{(j)}:=\prod_{p \in \mathcal{P}_{j}}(1-\alpha(p)), \quad \sum^{(j)}:=\sum_{\substack{d_{j} \mid P_{j} \\ \omega\left(d_{j}\right)=m_{j}+1}} \alpha\left(d_{j}\right), \tag{6.38}
\end{equation*}
$$

then

$$
\begin{aligned}
X \prod_{j=1}^{r} \sum_{\substack{d_{j} \mid P_{j} \\
\omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{j}\right) \alpha\left(d_{j}\right) & \leq X \prod_{j=1}^{r}\left(\prod^{(j)}+\sum^{(j)}\right) \\
& =X \prod_{p \in \mathcal{P}}(1-\alpha(p)) \prod_{j=1}^{r}\left(1+\sum^{(j)} / \prod^{(j)}\right)
\end{aligned}
$$

provided the division makes sense, i.e., provided $\alpha(p)<1$ for each $p \in \mathcal{P}$. Henceforth, we assume (as will be the case in all our applications) this condition on $\alpha$.

Recalling that $1+t \leq \exp (t)$, after estimating the remainder term of (6.37) trivially, we arrive at the following theorem:

Theorem 6.16 (Brun-Hooley sieve, upper bound). Let $\mathcal{P}=\dot{U}_{j=1}^{r} \mathcal{P}_{j}$ be a partition of $\mathcal{P}$. Suppose that $\alpha(p)<1$ for each $p \in \mathcal{P}$. For any choice of nonnegative even integers $m_{1}, \ldots, m_{r}$, we have

$$
\begin{align*}
S(\mathcal{A}, \mathcal{P}) \leq X \prod_{p \in \mathcal{P}}(1-\alpha(p)) \exp & \left(\sum_{j=1}^{r}\left(\sum^{(j)} / \prod^{(j)}\right)\right)  \tag{6.39}\\
& +O\left(\sum_{\substack{d_{1}, \ldots, d_{r} \\
d_{j} \mid P_{j}, \omega\left(d_{j}\right) \leq m_{j}}}\left|r\left(d_{1} \cdots d_{r}\right)\right|\right)
\end{align*}
$$

where $\prod^{(j)}$ and $\sum^{(j)}$ are defined, for $1 \leq j \leq r$, by (6.38), and the implied constant is absolute.
5.3. Applications of the upper bound. Define $R(N)$ as the number of (ordered) representations of $N$ as a sum of two primes, or equivalently, as the number of ordered prime pairs ( $p, N-p$ ). In Chapter 3, we conjectured that as $N \rightarrow \infty$ through even integers,

$$
R(N) \sim 2 C_{2} \frac{N}{(\log N)^{2}} \prod_{p \mid N, p>2} \frac{p-1}{p-2} .
$$

We now use the Brun-Hooley sieve to establish an upper bound for $R(N)$ of the conjecturally correct order of magnitude:

Theorem 6.17. For every even natural number $N$,

$$
R(N) \ll \frac{N}{(\log N)^{2}} \prod_{p \mid N}\left(1+\frac{1}{p}\right) .
$$

Let $N$ be an even natural number and define $\mathcal{A}:=\{n(N-n): 1 \leq n \leq$ $N\}$. Letting $\mathcal{P}$ be the set of all primes, we have for each choice of $z>0$,

$$
R(N) \leq 2 z+S(\mathcal{A}, \mathcal{P}, z)
$$

Indeed, if $N=n+(N-n)$ is a representation of $N$ as a sum of two primes, then either at least one of $n$ or $N-n$ lies in $[2, z]$ or both $n$ and $N-n$ have no prime factors $\leq z$. The former case occurs for no more than $2 z$ values of $n$, and the $n$ for which the latter holds (which necessarily satisfy $2 \leq n \leq N-2)$ are counted by $S(\mathcal{A}, \mathcal{P}, z)$.

We now choose our sifting parameters: Let $X=N$, and let $\alpha(d)=$ $\nu(d) / d$, where

$$
\nu(d):=\#\{n \bmod d: n(N-n) \equiv 0(\bmod d)\} ;
$$

then

$$
\alpha(p)= \begin{cases}1 / p & \text { if } p \mid N  \tag{6.40}\\ 2 / p & \text { if } p \nmid N .\end{cases}
$$

Because $N$ is even, $\alpha(p)<1$ for every prime $p$. Moreover,

$$
\begin{equation*}
A_{d}=X \alpha(d)+r(d) \quad \text { where } \quad|r(d)| \leq \nu(d) \quad \text { for all } d \mid P(z) . \tag{6.41}
\end{equation*}
$$

We think of $X=N$ as heading off towards infinity while $u>1$ is fixed. Our immediate goal is to show that if $u$ is fixed large enough, then

$$
S(\mathcal{A}, \mathcal{P}, z) \ll X \prod_{p \leq z}(1-\alpha(p)) \quad(X \rightarrow \infty), \quad \text { where } z:=X^{1 / u}
$$

To apply the Brun-Hooley sieve to this situation we need a partition of $\mathcal{P} \cap[2, z]$. We introduce the notation

$$
\eta=\log \log X
$$

and the choice of parameters

$$
\begin{equation*}
K:=1.57, \quad K_{1}:=1.571 \tag{6.42}
\end{equation*}
$$

For the present discussion it is only important that $1<K<K_{1}$, but this choice will be particularly effective for the lower bound applications of $\S 5.5$.

For large $X$, we have $\eta<z=X^{1 / u}$, so that if we define $R$ as the minimal integer with

$$
z^{1 / K^{R}}<\eta
$$

then $R \geq 1$. (Indeed, $R \rightarrow \infty$ with $X$.) For such $X$, we define

$$
z_{j}= \begin{cases}z^{1 / K^{j}} & \text { for } 0 \leq j \leq R-1 \\ \eta & \text { for } j=R \\ 1 & \text { for } j=R+1\end{cases}
$$

We partition $\mathcal{P} \cap[2, z]$ into the $r:=R+1$ sets

$$
\mathcal{P}_{j}:=\left\{p \in \mathcal{P}: z_{j}<p \leq z_{j-1}\right\} \quad(1 \leq j \leq R+1),
$$

and we define the corresponding nonnegative even integers $m_{1}, \ldots, m_{R+1}$ by putting

$$
m_{j}=2 j \quad(j=1, \ldots, R) \quad \text { and } \quad m_{R+1}=\infty ;
$$

here " $\infty$ " indicates that $m_{R+1}$ is chosen at least as large as the cardinality of $\mathcal{P}_{R+1}$. For definiteness, we take $m_{R+1}$ as the smallest even integer with this property. With this choice of $m_{R+1}$, the condition on a divisor $d$ of $P_{R+1}$ that it has no more than $m_{R+1}$ prime divisors becomes vacuous.

We are finally in a position to apply the upper bound (6.39) to our problem. By our choice of $m_{R+1}$,

$$
\begin{equation*}
\sum^{(R+1)}=\sum_{\substack{d_{R+1} \mid P_{R+1} \\ \omega\left(d_{R+1}\right)=m_{R+1}+1}} \alpha\left(d_{R+1}\right)=0 . \tag{6.43}
\end{equation*}
$$

Hence $\sum^{(j)} / \prod^{(j)}$ vanishes at $j=R+1$, and to estimate the main term of (6.39) it suffices to estimate the ratio $\sum^{(j)} / \prod^{(j)}$ for $j=1, \ldots, R$. The denominator is handled by the following lemma:

Lemma 6.18. As $x \rightarrow \infty$, we have

$$
\prod_{x<p \leq y}\left(1-\frac{2}{p}\right)=\frac{(\log x)^{2}}{(\log y)^{2}}\left(1+O\left(\frac{1}{\log x}\right)\right)
$$

uniformly for $y \geq x$.

Proof. Suppose $x \geq 4$; then $2 / p \leq 1 / 2$ for each $p \geq x$, so that $\log (1-2 / p)=$ $-2 / p+O\left((-2 / p)^{2}\right)$ with an absolute implied constant, and

$$
\begin{aligned}
\sum_{x<p \leq y} \log \left(1-\frac{2}{p}\right) & =-2 \sum_{x<p \leq y} \frac{1}{p}+O\left(\sum_{x<p \leq y} \frac{1}{p^{2}}\right) \\
& =-2\left(\log \frac{\log y}{\log x}+O\left(\frac{1}{\log x}\right)\right)+O\left(\frac{1}{x}\right) \\
& =\log \frac{(\log x)^{2}}{(\log y)^{2}}+O\left(\frac{1}{\log x}\right) .
\end{aligned}
$$

Exponentiating gives the result.

As $X \rightarrow \infty$, so do each of $z_{1}, \ldots, z_{R}$ (since each is at least $\eta$ ). Consequently, Lemma 6.18 implies that for large $X$ (and each $j=1,2, \ldots, R$ ),

$$
\begin{align*}
& \prod^{(j)}=\prod_{z_{j}<p \leq z_{j-1}}(1-\alpha(p)) \geq \prod_{z_{j}<p \leq z_{j-1}}\left(1-\frac{2}{p}\right) \\
& .44)=\frac{\left(\log z_{j}\right)^{2}}{\left(\log z_{j-1}\right)^{2}}\left(1+O\left(\frac{1}{\log z_{j}}\right)\right) \geq \frac{1}{K^{2}}\left(1+O\left(\frac{1}{\log \eta}\right)\right) \geq \frac{1}{K_{1}^{2}} . \tag{6.44}
\end{align*}
$$

Moreover, for $1 \leq j \leq R$, we have

$$
\begin{align*}
\sum^{(j)}=\sum_{\substack{d_{j} \mid P_{j} \\
\omega\left(d_{j}\right)=m_{j}+1}} \alpha\left(d_{j}\right) & \leq \frac{1}{\left(m_{j}+1\right)!}\left(\sum_{p \in \mathcal{P}_{j}} \alpha(p)\right)^{m_{j}+1} \\
& \leq \frac{1}{\left(m_{j}+1\right)!}\left(\sum_{p \in \mathcal{P}_{j}} \frac{2}{p}\right)^{m_{j}+1} \leq \frac{\left(2 \log K_{1}\right)^{m_{j}+1}}{\left(m_{j}+1\right)!} \tag{6.45}
\end{align*}
$$

provided $X$ is large enough, since in that case
$\sum_{z_{j}<p \leq z_{j-1}} \frac{2}{p}=2 \log \frac{\log z_{j-1}}{\log z_{j}}+O\left(\frac{1}{\log z_{j}}\right) \leq 2 \log K+O\left(\frac{1}{\log \eta}\right) \leq 2 \log K_{1}$.
Putting (6.44) and (6.45) together and recalling (6.43), we find that for large $X$,

$$
\sum_{j=1}^{R+1}\left(\sum^{(j)} / \prod^{(j)}\right) \leq K_{1}^{2} \sum_{j=1}^{R} \frac{\left(2 \log K_{1}\right)^{2 j+1}}{(2 j+1)!} \leq K_{1}^{2} \exp \left(2 \log K_{1}\right)
$$

This shows that the main term of (6.39) is bounded above by a constant multiple of $X \prod_{p \leq z}(1-\alpha(p))$. For any fixed $u>1$,

$$
\begin{equation*}
X \prod_{p \leq X^{1 / u}}(1-\alpha(p)) \geq \frac{1}{2} X \prod_{2<p \leq X^{1 / u}}(1-2 / p) \asymp X /(\log X)^{2} \quad(X \rightarrow \infty), \tag{6.46}
\end{equation*}
$$

so that to obtain the estimate $S(\mathcal{A}, \mathcal{P}, z) \ll X \prod_{p \leq z}(1-\alpha(p))$ we need only ensure that the sum appearing in the expression for the remainder term,

$$
\begin{equation*}
\sum_{\substack{d_{1}, \ldots, d_{R+1} \\ d_{j} \mid P_{j}, \omega\left(d_{j}\right) \leq m_{j}}}\left|r\left(d_{1} \cdots d_{R+1}\right)\right|, \tag{6.47}
\end{equation*}
$$

is of smaller order than $X /(\log X)^{2}$. We will show that for an appropriate choice of $u$, this sum is $\ll X^{\delta}$ for a constant $\delta<1$.

Observe that any product $d_{1} \cdots d_{R+1}$ appearing as an argument of $r(\cdot)$ in the sum (6.47) satisfies

$$
\begin{aligned}
d_{1} \cdots d_{R+1} & \leq\left(\prod_{j=1}^{R} z_{j-1}^{m_{j}}\right) \eta^{\eta} \\
& =X^{\frac{1}{u}\left(\sum_{j=1}^{R} m_{j} / K^{j-1}\right)} X^{\log \log X \log \log \log X / \log X} .
\end{aligned}
$$

Also,

$$
\sum_{j=1}^{R} \frac{m_{j}}{K^{j-1}} \leq \sum_{j=1}^{\infty} \frac{2 j}{K^{j-1}}=\frac{2 K^{2}}{(K-1)^{2}}=15.173 \ldots
$$

We fix a choice of $u$ exceeding $15.173 \ldots$, say $u=16$ for definiteness. Then for large enough $X$, we have $d_{1} \cdots d_{R+1} \leq X^{15.2 / 16}$ for every such product $d_{1} \cdots d_{R+1}$. For every $d$ dividing $P(z)$,

$$
|r(d)| \leq \nu(d)=\prod_{p \mid d} \nu(p) \leq 2^{\omega(d)} \leq \tau(d)
$$

Since each integer admits at most one representation in the form $d_{1} \cdots d_{R+1}$ (since the $d_{i}$ are supported on disjoint sets of primes), the sum (6.47) above is bounded by

$$
\sum_{n \leq X^{15.2 / 16}} \tau(n)=\sum_{n \leq X^{15.2 / 16}} \sum_{e \mid n} 1 \leq X^{15.2 / 16} \sum_{e \leq X^{15.2 / 16}} \frac{1}{e} \ll X^{15.2 / 16} \log X
$$

It follows that for all large $X$,

$$
\begin{aligned}
S\left(\mathcal{A}, \mathcal{P}, X^{\frac{1}{16}}\right) & \ll X \prod_{p \leq X \frac{1}{16}}(1-\alpha(p)) \\
& =X \prod_{\substack{p \leq X \frac{1}{16} \\
p \nmid N}}\left(1-\frac{2}{p}\right) \prod_{\substack{\left.p \leq X \frac{1}{p} 16 \\
p \right\rvert\, N}}\left(1-\frac{1}{p}\right) .
\end{aligned}
$$

Since $(1-2 / p) \leq(1-1 / p)^{2}$, we find that

$$
\begin{aligned}
S\left(\mathcal{A}, \mathcal{P}, X^{\frac{1}{16}}\right) & \leq X \prod_{\substack{p \leq X \frac{1}{16} \\
p \nmid N}}\left(1-\frac{1}{p}\right)^{2} \prod_{\substack{\left.p \leq X \frac{1}{p-16} \\
p \right\rvert\, N}}\left(1-\frac{1}{p}\right) \\
& =X \prod_{p \leq X \frac{1}{16}}\left(1-\frac{1}{p}\right)^{2} \prod_{\substack{\left.p \leq X \frac{1}{16} \\
p \right\rvert\, N}}\left(1-\frac{1}{p}\right)^{-1} \\
& \ll \frac{X}{(\log X)^{2}} \prod_{p \mid N}\left(1-\frac{1}{p}\right)^{-1}
\end{aligned}
$$

Noting that

$$
\prod_{p \mid N}\left(1-\frac{1}{p}\right)^{-1} / \prod_{p \mid N}\left(1+\frac{1}{p}\right)=\prod_{p \mid N}\left(1-\frac{1}{p^{2}}\right)^{-1} \leq \sum_{n=1}^{\infty} \frac{1}{n^{2}}<\infty
$$

we conclude that for large $X$,

$$
\begin{equation*}
S\left(\mathcal{A}, \mathcal{P}, X^{1 / 16}\right) \ll \frac{X}{(\log X)^{2}} \prod_{p \mid N}\left(1+\frac{1}{p}\right) \tag{6.48}
\end{equation*}
$$

Consequently, for all large positive even numbers $N$,

$$
\begin{aligned}
R(N) & \leq S\left(\mathcal{A}, \mathcal{P}, X^{1 / 16}\right)+2 X^{1 / 16} \\
& \ll \frac{X}{(\log X)^{2}} \prod_{p \mid N}\left(1+\frac{1}{p}\right)=\frac{N}{(\log N)^{2}} \prod_{p \mid N}\left(1+\frac{1}{p}\right) .
\end{aligned}
$$

This gives the assertion of Theorem 6.17 for sufficiently large $N$, but for bounded $N$ the theorem is trivial.

The proof we have given applies mutatis mutandis to the generalized prime twin problem, i.e., the problem of estimating

$$
\pi_{N}(x):=\#\{p \leq x: p, p+N \text { are both prime }\} .
$$

Indeed, let $N$ be a positive even integer, and define the sequence

$$
\mathcal{A}:=\{n(n+N): 1 \leq n \leq x\} .
$$

Then

$$
\pi_{N}(x) \leq z+S(\mathcal{A}, \mathcal{P}, z)
$$

for any choice of positive $z$. To estimate $S(\mathcal{A}, \mathcal{P}, z)$, we take $X=x$ and choose $\alpha(d)=\nu(d) / d$, where $\nu(d)$ is the number of solutions to the congruence $n(N+n) \equiv 0(\bmod d)$. Then $\alpha(d)$ is again given by (6.40). If we now choose $z$, the $z_{j}$, the partition $\mathcal{P}_{j}$, and the $m_{j}$ exactly as before, the same proof as above shows that (6.48) holds for all sufficiently large $X$, say $X \geq x_{0}$. Moreover, both $x_{0}$ and the implied constant in (6.48) are independent of $N$. So, for $x \geq x_{0}$,

$$
\begin{aligned}
\pi_{N}(x) & \ll X^{1 / 16}+S\left(\mathcal{A}, \mathcal{P}, X^{1 / 16}\right) \\
& \ll x^{1 / 16}+\frac{x}{(\log x)^{2}} \prod_{p \mid N}\left(1+\frac{1}{p}\right) \ll \frac{x}{(\log x)^{2}} \prod_{p \mid N}\left(1+\frac{1}{p}\right),
\end{aligned}
$$

uniformly in $N$. Since $\pi_{N}(x)$ is trivially bounded by $x_{0}$ for $2 \leq x \leq x_{0}$, the same upper estimate for $\pi_{N}(x)$ remains valid for all $x \geq 2$ and all even natural numbers $N$ (with perhaps a different implied constant). So we have proved:

Theorem 6.19. Let $N$ be a positive even integer. Then for $x \geq 2$,

$$
\pi_{N}(x) \ll \frac{x}{(\log x)^{2}} \prod_{p \mid N}\left(1+\frac{1}{p}\right),
$$

where the implied constant is absolute.
5.4. The lower bound. We turn now to the problem of bounding $S(\mathcal{A}, \mathcal{P})$ from below. A natural temptation here is to simply parallel what we did in the upper bound case: If we suppose $m_{1}, \ldots, m_{r}$ to be $r$ odd natural numbers, then for each $j$,

$$
\sum_{\substack{d_{j}\left|n, d_{j}\right| P_{j} \\ \omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{j}\right) \leq \sum_{d_{j}\left|n, d_{j}\right| P_{j}} \mu\left(d_{j}\right) .
$$

But since it is (generally) not the case that for every $1 \leq j \leq r$, both sides of this inequality are nonnegative, we cannot simply take the product of both sides over $j$ and expect the inequality to be preserved.

So we require a different approach. By Lemma 6.15 (with $\mathcal{P}, P$ replaced by $\mathcal{P}_{j}, P_{j}$ ), for any choice of nonnegative even integers $m_{1}, \ldots, m_{r}$, we have

$$
\begin{equation*}
0 \leq \sum_{\substack{d_{j}\left|n, d_{j}\right| P_{j} \\ \omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{j}\right)-\sum_{d_{j}\left|n, d_{j}\right| P} \mu\left(d_{j}\right) \leq \sum_{\substack{d_{j}\left|n, d_{j}\right| P \\ \omega\left(d_{j}\right)=m_{j}+1}} 1 \quad(1 \leq j \leq r) \tag{6.49}
\end{equation*}
$$

These bounds allow us to a coax a lower bound for the sifting function

$$
\begin{equation*}
s(n)=\prod_{j=1}^{r} \sum_{d_{j}\left|n, d_{j}\right| P_{j}} \mu\left(d_{j}\right) \tag{6.50}
\end{equation*}
$$

out of the following general inequality:
Lemma 6.20 ([FH00, Lemma 1]). Suppose that $0 \leq x_{j} \leq y_{j}$ for $1 \leq j \leq r$. Then

$$
x_{1} \cdots x_{r} \geq y_{1} \cdots y_{r}-\sum_{l=1}^{r}\left(y_{l}-x_{l}\right) \prod_{\substack{j=1 \\ j \neq l}}^{r} y_{j} .
$$

Proof. The result holds with equality when $r=1$. If the lemma holds for $r-1$ for a certain $r \geq 2$, then

$$
\begin{aligned}
y_{1} \cdots y_{r}-x_{1} \cdots x_{r} & =\left(y_{1} \cdots y_{r-1}-x_{1} \cdots x_{r-1}\right) y_{r}+\left(x_{1} \cdots x_{r-1}\right)\left(y_{r}-x_{r}\right) \\
& \leq\left(y_{1} \cdots y_{r-1}-x_{1} \cdots x_{r-1}\right) y_{r}+\left(y_{1} \cdots y_{r-1}\right)\left(y_{r}-x_{r}\right) \\
& \leq \sum_{l=1}^{r-1}\left(y_{l}-x_{l}\right) \prod_{\substack{j=1 \\
j \neq l}}^{r} y_{j}+\left(y_{r}-x_{r}\right) \prod_{\substack{j=1 \\
j \neq r}}^{r} y_{j},
\end{aligned}
$$

which is just $\sum_{l=1}^{r}\left(y_{l}-x_{l}\right) \prod_{\substack{j=1 \\ j \neq l}}^{r} y_{j}$. So the result follows by induction.

Assuming $m_{1}, \ldots, m_{r}$ are nonnegative even integers, we apply Lemma 6.20 with

$$
x_{j}:=\sum_{d_{j}\left|n, d_{j}\right| P_{j}} \mu\left(d_{j}\right), \quad y_{j}:=\sum_{\substack{d_{j}\left|n, d_{j}\right| P_{j} \\ \omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{j}\right) .
$$

Equation (6.49) implies that the hypotheses of Lemma 6.20 are satisfied and gives us an upper bound on the terms $y_{l}-x_{l}$. Using this bound in Lemma 6.20 and recalling (6.50), we obtain

$$
s(n) \geq \prod_{j=1}^{r} \sum_{\substack{d_{j}\left|n, d_{j}\right| P_{j} \\ \omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{j}\right)-\sum_{l=1}^{r}\left(\sum_{\substack{d_{l}\left|n, d_{l}\right| P_{l} \\ \omega\left(d_{l}\right)=m_{l}+1}} 1\right) \prod_{\substack{j=1 \\ j \neq l}}^{r}\left(\sum_{\substack{d_{j}\left|n, d_{j}\right| P_{j} \\ \omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{j}\right)\right) .
$$

Summing over $n \in \mathcal{A}$ shows that

$$
\begin{align*}
& S(\mathcal{A}, \mathcal{P}) \geq \sum_{\substack{d_{1}, \ldots, d_{r} \\
d_{j} \mid P_{j}, \omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{1}\right) \cdots \mu\left(d_{r}\right) A_{d_{1} \cdots d_{r}}  \tag{6.51}\\
&\left.-\sum_{l=1}^{r} \sum_{\substack{\left.d_{1}, \ldots, d_{r} \\
d_{j}\left|P_{j}, \omega\left(d_{j}\right) \leq m_{j} \\
d_{l}\right| P_{l}, \omega \neq l\right)}} \frac{\mu\left(d_{l}\right)=m_{l}+1}{}\right) \cdots \mu\left(d_{r}\right) \\
& \mu\left(d_{l}\right) d_{d_{1} \cdots d_{r}} .
\end{align*}
$$

Writing $A_{d}=X \alpha(d)+r(d)$, the right-hand side of (6.51) becomes

$$
\begin{equation*}
X \prod_{j=1}^{r} \sum_{\substack{d_{j} \mid P_{j} \\ \omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{j}\right) \alpha\left(d_{j}\right)-X \sum_{l=1}^{r} \sum_{\substack{d_{l} \mid P_{l} \\ \omega\left(d_{l}\right)=m_{l}+1}} \alpha\left(d_{l}\right) \prod_{j \neq l} \sum_{\substack{d_{j} \mid P_{j} \\ \omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{j}\right) \alpha\left(d_{j}\right), \tag{6.52}
\end{equation*}
$$

up to an error term that is (with an absolute implied constant)

$$
\ll \sum_{\substack{d_{j} \mid P_{j}(1 \leq j \leq r) \\ \theta_{d_{1}}, \ldots, d_{r}}}\left|r\left(d_{1} \cdots d_{r}\right)\right| .
$$

Here $\theta_{d_{1}, \ldots, d_{r}}$ denotes the condition that there exist $r-1$ indices $j, 1 \leq j \leq r$, for which $\omega\left(d_{j}\right) \leq m_{j}$, while the remaining index satisfies $\omega\left(d_{j}\right) \leq m_{j}+1$.

Assume, as we did for the upper bound, that $\alpha(p)<1$ for each $p \in \mathcal{P}$. Lemma 6.6 implies that for each $1 \leq j \leq r$,

$$
\sum_{\substack{d_{j} \mid P_{j} \\ \omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{j}\right) \alpha\left(d_{j}\right) \geq \prod_{p \in \mathcal{P}_{j}}(1-\alpha(p))>0,
$$

so that the main term in (6.52) is

$$
\begin{aligned}
& X\left(1-\sum_{1 \leq l \leq r} \frac{\sum_{d_{l} \mid P_{l}, \omega\left(d_{l}\right)=m_{l}+1} \alpha\left(d_{l}\right)}{\sum_{d_{l} \mid P_{l}, \omega\left(d_{l}\right) \leq m_{l}} \mu\left(d_{l}\right) \alpha\left(d_{l}\right)}\right) \prod_{j=1}^{r} \sum_{\substack{d_{j} \mid P_{j} \\
\omega\left(d_{j}\right) \leq m_{j}}} \mu\left(d_{j}\right) \alpha\left(d_{j}\right) \\
& \quad \geq X \prod_{p \in \mathcal{P}}(1-\alpha(p))\left(1-\sum_{1 \leq l \leq r}\left(\sum_{\substack{d_{l} \mid P_{l} \\
\omega\left(d_{l}\right)=m_{l}+1}} \alpha\left(d_{l}\right) / \prod_{p \in \mathcal{P}_{l}}(1-\alpha(p))\right)\right) .
\end{aligned}
$$

Summarizing, we have proved the following theorem:
Theorem 6.21 (Brun-Hooley sieve, lower bound). Let $\mathcal{P}=\dot{\bigcup}_{j=1}^{r} \mathcal{P}_{j}$ be a partition of $\mathcal{P}$. Suppose that $\alpha(p)<1$ for each $p \in \mathcal{P}$. For any choice of nonnegative even integers $m_{1}, \ldots, m_{r}$, we have

$$
\begin{aligned}
S(\mathcal{A}, \mathcal{P}) \geq X \prod_{p \in \mathcal{P}}(1-\alpha(p))\left(1-\sum_{j=1}^{r}\right. & \left.\left(\sum^{(j)} / \prod^{(j)}\right)\right) \\
& +O\left(\sum_{\substack{d_{j} \mid P_{j}(1 \leq j \leq r) \\
\theta_{d_{1}}, \ldots, d_{r}}}\left|r\left(d_{1} \cdots d_{r}\right)\right|\right),
\end{aligned}
$$

where $\prod^{(j)}$ and $\sum^{(j)}$ are defined, for $1 \leq j \leq r$, by (6.38), and the implied constant is absolute.
5.5. Applications of the lower bound. We now prove the two remarkable theorems of Brun mentioned in the introduction: Every large even integer is a sum of two 9 -almost primes, and there exist infinitely many pairs of 9 -almost primes differing by 2 .

Our setup for attacking these problems is the same as that used in attacking the analogous upper bound problems considered in $\S 5.3$. For the first of these, we assume $N$ is an even natural number, and we take $\mathcal{A}:=$ $\{n(N-n): 1 \leq n \leq N\}$. As before, we let $\mathcal{P}$ be the set of all primes.

Suppose that we have a positive even integer $N$ and a $u>1$ for which

$$
\begin{equation*}
S\left(\mathcal{A}, \mathcal{P}, N^{1 / u}\right)>0 \tag{6.53}
\end{equation*}
$$

Then there exists an $n, 1 \leq n \leq N$, such that both $n$ and $N-n$ have all their prime divisors exceeding $N^{1 / u}$; since both $n$ and $N-n$ are bounded by $N$, each must have at most $u$ prime divisors. We will show that if we choose $u$ large enough, (6.53) holds for all sufficiently large $N$ (depending on $u$ ). Brun's results then follow from a quantitative determination of which $u$ are "large enough".

For the most part, we may choose our sieving parameters as in $\S 5.3$, so that $X=N$ and $\alpha$ is given by (6.40). With $u$ a parameter to be chosen later, we define the partition of $\mathcal{P} \cap[2, z]$ into sets $\mathcal{P}_{j}$ as in $\S 5.3$. However, the choice of the corresponding $m_{j}$ requires more care.

To describe this choice, suppose for the moment that we have constructed a sequence $\left\{n_{i}\right\}_{i=1}^{\infty}$ of nonnegative even integers satisfying the two inequalities

$$
\begin{align*}
& \sum_{j=1}^{\infty} \frac{\left(2 \log K_{1}\right)^{n_{j}+1}}{\left(n_{j}+1\right)!}<\frac{1}{K_{1}^{2}},  \tag{6.54}\\
& \Gamma:=1+\sum_{j=1}^{\infty} \frac{n_{j}}{K^{j-1}}<\infty, \tag{6.55}
\end{align*}
$$

where $K$ and $K_{1}$ are given by (6.42). We fix $u>\Gamma$ and define (with same meaning of " $\infty$ " as in §5.3)

$$
m_{j}=n_{j} \quad(1 \leq j \leq R), \quad m_{R+1}=\infty .
$$

Then for all large $X$, we have (recalling (6.43), (6.44), (6.45))

$$
\begin{aligned}
\sum_{j=1}^{R+1}\left(\sum^{(j)} / \prod^{(j)}\right) & =\sum_{j=1}^{R}\left(\sum^{(j)} / \prod^{(j)}\right) \\
& \leq K_{1}^{2} \sum_{j=1}^{R} \sum^{(j)} \leq K_{1}^{2} \sum_{j=1}^{R} \frac{\left(2 \log K_{1}\right)^{m_{j}+1}}{\left(m_{j}+1\right)!} \leq 1-\epsilon
\end{aligned}
$$

for a positive constant $\epsilon$, by (6.54). This implies that the main term in the lower bound

$$
\begin{align*}
S(\mathcal{A}, \mathcal{P}) \geq X \prod_{p \in \mathcal{P}}(1-\alpha(p)) & \left(1-\sum_{1 \leq j \leq R+1}\left(\sum^{(j)} / \prod^{(j)}\right)\right)  \tag{6.56}\\
& +O\left(\sum_{\substack{1 \leq\left|P_{j}\right| P_{j}(1 \leq j \leq R+1) \\
\theta_{d_{1}}, \ldots, d_{R+1}}}\left|r\left(d_{1} \cdots d_{R+1}\right)\right|\right),
\end{align*}
$$

is (cf. (6.46))

$$
\gg X \prod_{p \leq X^{1 / u}}(1-\alpha(p)) \gg X /(\log X)^{2} \quad(X \rightarrow \infty) .
$$

The $O$-term can be treated much as in $\S 5.3$ : The largest value of $d_{1} \cdots d_{R+1}$ appearing as an argument of $r(\cdot)$ is bounded above by

$$
X^{\frac{1}{u}\left(1+\sum_{j=1}^{R} m_{j} / K^{j-1}\right)} X^{\log \log X \log \log \log X / \log X} \leq X^{\Gamma / u+o(1)} \leq X^{\delta}
$$

for all large $X$, where $\delta:=\frac{1}{2}(1+\Gamma / u)$. Notice that $\delta<1$. The argument of $\S 5.3$ shows that the $O$-term in (6.56) is $\ll X^{\delta} \log X$, which is $o\left(X /(\log X)^{2}\right)$. So with this choice of parameters, we obtain (6.53) in the stronger form

$$
S\left(\mathcal{A}, \mathcal{P}, X^{1 / u}\right) \gg X /(\log X)^{2} \quad(X \rightarrow \infty)
$$

It remains to construct a suitable sequence $\left\{n_{i}\right\}$. It is not hard to see that (6.54) and (6.55) will be satisfied with the simple choice $n_{i}=b+2(i-1)$ ( $i \geq 1$ ), if we pick $b$ to be a suitably large even natural number. However, this construction leads to an unnecessarily bloated value of $\Gamma$, so that while we still obtain a statement of the form "every large even $N$ is a sum of two numbers with $O(1)$ prime factors", the $O(1)$ term dictating the number of summands is larger than we might like. We do better if we use the greedy algorithm to pick the first several $n_{i}$ (which play the largest role in determining the size of $\Gamma$ ): Choose as many of the initial $n_{i}$ to be 2 as (6.54) allows, then as many of the subsequent $n_{i}$ to be 4 as allowed, etc.

Using a calculator or computer, we find that the sequence obtained in this way begins

$$
n_{1}=n_{2}=n_{3}=2, \quad n_{4}=\cdots=n_{10}=4, \quad n_{11}=\cdots=n_{24}=6 .
$$

Instead of continuing in this manner, we make the simple choice

$$
n_{25}=8+2(j-25) \quad(j \geq 25)
$$

Then, setting $L:=2 \log K_{1}$,

$$
\begin{aligned}
\frac{1}{K_{1}^{2}} & -\sum_{j=1}^{\infty} \frac{\left(2 \log K_{1}\right)^{n_{j}+1}}{\left(n_{j}+1\right)!} \\
& \geq \frac{1}{K_{1}^{2}}-\sum_{j=1}^{3} \frac{L^{3}}{3!}-\sum_{j=4}^{10} \frac{L^{5}}{5!}-\sum_{j=11}^{24} \frac{L^{7}}{7!}-\sum_{j=25}^{\infty} \frac{L^{9+2(j-25)}}{(9+2(j-25))!} \\
& \geq \frac{1}{K_{1}^{2}}-3 \frac{L^{3}}{3!}-7 \frac{L^{5}}{5!}-14 \frac{L^{7}}{7!}-\frac{L^{9} / 9!}{1-L^{2} /(11 \cdot 10)}=0.00003 \ldots>0,
\end{aligned}
$$

so that (6.54) holds in this case. Also,

$$
\begin{aligned}
\Gamma & =1+\sum_{j=1}^{3} \frac{2}{K^{j-1}}+\sum_{j=4}^{10} \frac{4}{K^{j-1}}+\sum_{j=11}^{24} \frac{6}{K^{j-1}}+\sum_{j=25}^{\infty} \frac{8+2(j-25)}{K^{j-1}} \\
& =1+\sum_{j=1}^{3} \frac{2}{K^{j-1}}+\sum_{j=4}^{10} \frac{4}{K^{j-1}}+\sum_{j=11}^{24} \frac{6}{K^{j-1}}+\frac{2(4 K-3)}{K^{23}(K-1)^{2}}=7.993 \ldots
\end{aligned}
$$

Thus (6.55) holds. Moreover, we can take $u=7.995$, say. Doing so, we obtain an even stronger theorem than that stated in the introduction: Every large enough even $N$ may be represented as a sum of two natural numbers
each of which has no more than 7 prime divisors, and the number of such representations is $\gg X /(\log X)^{2}=N /(\log N)^{2}$ as $N \rightarrow \infty$.

In like manner, one can show that there are $\gg x /(\log x)^{2}$ positive integers $n \leq x$ for which both $n$ and $n+N$ have no prime divisor $\leq x^{1 / 7.995}$, uniformly in the choice of the even natural number $N$. Suppose now that $N$ is fixed; then for large enough $x$, we have

$$
n \leq n+N \leq x+N<\left(x^{1 / 7.995}\right)^{8} ;
$$

it follows that there are

$$
>_{N} x /(\log x)^{2} \quad(x \rightarrow \infty)
$$

integers $n \leq x$ for which both $n$ and $n+N$ have no more than 7 prime divisors. When $N=2$ we obtain Brun's statement (with 9 replaced by the superior constant 7).

Note that $K$ and $K_{1}$ in (6.42) were chosen to minimize the quantity $\Gamma$, which is the limiting factor in how small we are allowed to select $u$. Their numerical values were found by computer (cf. [FH00, pp. 347-348]).

## 6. An application to the Goldbach problem

While sieve methods are now part of the standard tool chest of analytic number theory, this was not always the case. In the monograph of Halberstam \& Richert [HR74, p. 6], the story is told of how Landau left Brun's manuscript untouched in a drawer for six years until hearing of a striking application made by the Russian mathematician Schnirelmann [Sch33]:
Theorem 6.22. There is an absolute constant $S$ with the following property: Every integer $n>1$ can be written as a sum of at most $S$ prime numbers.

Our objective in this section is to prove Theorem 6.22.
6.1. Schnirelmann density. Write $\mathbf{N}_{0}$ for the set of nonnegative integers. In what follows we use script letters to denote subsets of $\mathbf{N}_{0}$ and use the corresponding Roman letters for their counting functions. Even though such sets may contain zero, it is convenient to define our counting functions so that only positive elements are tallied; thus, e.g.,

$$
A(n)=\#\{a \in \mathcal{A}: 1 \leq a \leq n\} .
$$

If $\mathcal{A}, \mathcal{B} \subset \mathbf{N}$, we define the sumset $\mathcal{A} \oplus \mathcal{B}$ by

$$
\mathcal{A} \oplus \mathcal{B}:=\{a+b: a \in \mathcal{A}, b \in \mathcal{B}\} .
$$

For $h \in \mathbf{N}$, we put

$$
h \mathcal{A}:=\overbrace{\mathcal{A} \oplus \cdots \oplus \mathcal{A}}^{h \text { summands }} .
$$

We say that $\mathcal{A}$ is a basis of finite order if $h \mathcal{A}=\mathbf{N}_{0}$ for some $h \in \mathbf{N}$. In this case the smallest such $h$ is called the order of the basis. For example, if $\mathcal{A}=\left\{n^{2}: n \in \mathbf{Z}\right\}$, then $\mathcal{A}$ is a basis of order 4. In fact, if $k$ is any integer with $k \geq 2$, then $\left\{n^{k}: n \in \mathbf{N}_{0}\right\}$ is a basis of finite order by the Hilbert-Waring Theorem considered in Chapter 5.

For each subset $A \subset \mathbf{N}_{0}$, we define the Schnirelmann density $\delta(\mathcal{A})$ of $\mathcal{A}$ by

$$
\delta(\mathcal{A}):=\inf _{n=1,2,3, \ldots} \frac{A(n)}{n} .
$$

This definition is a bit odd; unlike (e.g.) the notion of asymptotic density, the presence (or absence) of small numbers in $\mathcal{A}$ has a disproportionate impact. The most extreme instance of this is that $\mathcal{A}$ automatically has Schnirelmann density zero whenever $1 \notin \mathcal{A}$. Moreover, the only way that a set $\mathcal{A}$ can have Schnirelmann density 1 is if $\mathcal{A}$ contains every natural number. Despite these peculiarities, the Schnirelmann density is a very convenient measure of size for questions in additive number theory. Indeed, Schnirelmann succeeded in proving the following very useful criterion for a set to be a basis of finite order:

Theorem 6.23 (Schnirelmann's basis theorem). Let $\mathcal{A}$ be a subset of $\mathbf{N}_{0}$ with $0 \in \mathcal{A}$ and $\delta(\mathcal{A})>0$. Then $\mathcal{A}$ is a basis of finite order.

The proof requires two simple lemmas.
Lemma 6.24. If $\mathcal{A}$ and $\mathcal{B}$ are sets of nonnegative integers, each containing 0 , and $\delta(\mathcal{A})+\delta(\mathcal{B}) \geq 1$, then $\mathcal{A} \oplus \mathcal{B}=\mathbf{N}_{0}$. In particular, if $0 \in \mathcal{A}$ and $\delta(\mathcal{A}) \geq 1 / 2$, then $2 \mathcal{A}=\mathbf{N}_{0}$.

Proof. We will show that each $n \in \mathbf{N}_{0}$ belongs to the sumset $\mathcal{A} \oplus \mathcal{B}$. Suppose that $a_{0}=0<a_{1}<a_{2}<\cdots$ is an enumeration of $\mathcal{A}$ and that $0=b_{0}<b_{1}<b_{2}<\cdots$ is an enumeration of $\mathcal{B}$. Let $n \in \mathbf{N}_{0}$, and consider the following list of nonnegative integers from $[0, n]$ :

$$
0=a_{0}, a_{1}, \ldots, a_{A(n)}, n=n-b_{0}, n-b_{1}, \ldots, n-b_{B(n)} .
$$

This list has length

$$
(A(n)+1)+(B(n)+1) \geq \delta(\mathcal{A}) n+\delta(\mathcal{B}) n+2 \geq n+2>n+1
$$

Since there are only $n+1$ integers in the interval $[0, n]$, it must be that for some pair of $i$ and $j$ with $0 \leq i \leq A(n)$ and $0 \leq j \leq B(n)$, we have $a_{i}=n-b_{j}$. But then $n=a_{i}+b_{j} \in \mathcal{A} \oplus \mathcal{B}$.

Lemma 6.25. If $\mathcal{A}$ and $\mathcal{B}$ are sets of nonnegative integers, each containing 0 , then $\delta(\mathcal{A} \oplus \mathcal{B}) \geq \delta(\mathcal{A})+\delta(\mathcal{B})-\delta(\mathcal{A}) \delta(\mathcal{B})$.

Proof. Let $n \in \mathbf{N}$, and let $0<a_{1}<a_{2}<\cdots<a_{A(n)} \leq n$ be a list of the elements of $\mathcal{A} \cap[1, n]$. Define intervals $I_{j}$ for $0 \leq j \leq A(n)$ by putting $I_{0}=\left(0, a_{1}\right), I_{1}=\left(a_{1}, a_{2}\right), I_{2}=\left(a_{2}, a_{3}\right), \ldots, I_{A(n)-1}=\left(a_{A(n)-1}, a_{A(n)}\right)$, and $I_{A(n)}=\left(a_{A(n)}, n\right]$. We now estimate $\#(\mathcal{A} \oplus \mathcal{B}) \cap I_{j}$ for each $j$.

For $j=0$, we have $\#(\mathcal{A} \oplus \mathcal{B}) \cap I_{0} \geq B\left(a_{1}-1\right)$, since if $b \in \mathcal{B} \cap\left[1, a_{1}-1\right]$, then $0+b \in(\mathcal{A} \oplus \mathcal{B}) \cap I_{0}$. Similarly, for $1 \leq j<A(n)$, we have $\#(\mathcal{A} \oplus \mathcal{B}) \cap I_{j} \geq$ $B\left(a_{j+1}-a_{j}-1\right)$, since if $b \in \mathcal{B} \cap\left[1, a_{j+1}-a_{j}-1\right]$, then $a_{j}+b \in(\mathcal{A} \oplus \mathcal{B}) \cap I_{j}$. Finally, $\#(\mathcal{A} \oplus \mathcal{B}) \cap I_{A(n)} \geq B\left(n-a_{A(n)}\right)$, since if $b \in \mathcal{B} \cap\left[1, n-a_{A(n)}\right]$, then $a_{A(n)}+b \in(\mathcal{A} \oplus \mathcal{B}) \cap I_{A(n)}$. Moreover, since $0 \in \mathcal{B}$, we know also that $\mathcal{A} \oplus \mathcal{B} \supset \mathcal{A}$. Hence,

$$
\begin{aligned}
(A \oplus B)(n) & \geq A(n)+\sum_{i=0}^{A(n)} \#(\mathcal{A} \oplus \mathcal{B}) \cap I_{j} \\
& \geq A(n)+B\left(a_{1}-1\right)+\sum_{i=1}^{A(n)-1} B\left(a_{i+1}-a_{i}-1\right)+B\left(n-a_{n}\right) .
\end{aligned}
$$

Since $B(m) \geq \delta(\mathcal{B}) m$ for each $m \in \mathbf{N}_{0}$, this is at least

$$
A(n)+\delta(\mathcal{B})\left(\left(a_{1}-1\right)+\sum_{i=1}^{A(n)-1}\left(a_{i+1}-a_{i}-1\right)+n-a_{A(n)}\right) .
$$

But $A(n) \geq \delta(\mathcal{A}) n$, so that

$$
\begin{aligned}
(A \oplus B)(n) & \geq \delta(\mathcal{A}) n(1-\delta(\mathcal{B}))+\delta(\mathcal{B}) n \\
& =n(\delta(\mathcal{A})+\delta(\mathcal{B})-\delta(\mathcal{A}) \delta(\mathcal{B})) .
\end{aligned}
$$

Since $n$ was arbitrary, the assertion of the lemma follows from the definition of Schnirelmann density.

Proof of Theorem 6.23. Taking $\mathcal{A}=\mathcal{B}$ in Lemma 6.25, we find $\delta(2 \mathcal{A}) \geq$ $2 \delta(\mathcal{A})-\delta(\mathcal{A})^{2}$. Said differently, $1-\delta(2 \mathcal{A}) \leq(1-\delta(\mathcal{A}))^{2}$. Starting from this inequality, an easy induction shows that for every $k \geq 1$,

$$
1-\delta\left(2^{k} \mathcal{A}\right) \leq(1-\delta(\mathcal{A}))^{2^{k}}
$$

Since $\delta(\mathcal{A})>0$, we can choose a natural number $k$ for which the righthand side of this inequality is at most $1 / 2$. Then $\delta\left(2^{k} \mathcal{A}\right) \geq 1 / 2$, and so $2^{k+1} \mathcal{A}=\mathbf{N}_{0}$ by Lemma 6.24. So $\mathcal{A}$ is a basis of order at most $2^{k+1}$.

Remark. A theorem of Mann [Man42], strengthening Lemma 6.25, asserts that if $\mathcal{A}$ and $\mathcal{B}$ are subsets of $\mathbf{N}_{0}$ with $0 \in \mathcal{A} \cap \mathcal{B}$, then $\delta(\mathcal{A} \oplus \mathcal{B}) \geq$ $\min \{1, \delta(\mathcal{A})+\delta(\mathcal{B})\}$. This had been conjectured by Landau \& Schnirelmann. An immediate consequence of Mann's theorem is that under the hypotheses of Theorem $6.23, \mathcal{A}$ is a basis of order at most $\lceil 1 / \delta(\mathcal{A})\rceil$. For a discussion of

Mann's theorem and subsequent related developments (including the important work of Kneser), see the volumes of Ostmann mentioned in the notes at the end of this chapter. There is also some discussion of these results in the appealing survey [PS95].
6.2. Proof of Theorem 6.22. Observe that if $\mathcal{A} \subset \mathbf{N}_{0}$ has positive lower density, in the sense that

$$
\begin{equation*}
\liminf _{x \rightarrow \infty} \frac{A(x)}{x}>0 \tag{6.57}
\end{equation*}
$$

then $\mathcal{B}:=\{0,1\} \cup \mathcal{A}$ has positive Schnirelmann density. Indeed, (6.57) implies that for some $\delta_{0}>0$ and $N_{0} \in \mathbf{N}$, we have $A(N) \geq \delta_{0} N$ for all $N \geq N_{0}$. But then $\delta(\mathcal{B}) \geq \min \left\{\delta_{0}, 1 / N_{0}\right\}>0$. Since also $0 \in \mathcal{B}$, we may apply Theorem 6.23 to deduce that $\mathcal{B}$ is a basis of finite order. We will shortly make use of these observations for an appropriately chosen set $\mathcal{A}$.

Recall that for a natural number $N$, the number of ordered representations of $N$ as a sum of two primes is denoted by $R(N)$. For each $N \geq 2$, we have

$$
\begin{equation*}
R(N) \ll \frac{N}{(\log N)^{2}} \prod_{p \mid N}\left(1+\frac{1}{p}\right) . \tag{6.58}
\end{equation*}
$$

(This was proved in $\S 5.3$ when $N$ is even. If $N$ is odd, then $R(N) \leq 2$ and so (6.58) is trivial.) We now let

$$
\mathcal{A}:=\{N \in \mathbf{N}: R(N)>0\} .
$$

We will prove the following:
Theorem 6.26. The set $\mathcal{A}$ has positive lower density.
Once this is proved, Theorem 6.23 follows easily. Indeed, let $\mathcal{B}=\mathcal{A} \cup$ $\{0,1\}$, so that from the above discussion $\mathcal{B}$ is a basis of finite order $h \geq 1$, say. Then for every integer $n \geq 2$, we can write

$$
n-2=p_{1}+p_{2}+\cdots+p_{2 k}+\overbrace{1+1+\cdots+1}^{l \text { summands }},
$$

say, where the $p_{i}$ are primes, $k$ and $l$ are nonnegative integers, and $k+l \leq h$. Then

$$
n=p_{1}+\cdots+p_{2 k}+(l+2)
$$

Since $l+2 \geq 2$, it can be written as a sum of 2 s and 3 s , where the number of summands is at most $(l+2) / 2 \leq h / 2+1$. This means that $n$ has a representation as a sum of at most $2 k+h / 2+1 \leq 5 h / 2+1$ primes. Theorem 6.23 follows with $S=5 h / 2+1$.

The main tool needed in the proof of Theorem 6.26 is the upper bound (6.58). It is initially surprising that an upper bound for $R(N)$ would be of
use in establishing a lower density result. But this seeming paradox is easily explained: As we will see shortly, it is a simple matter to obtain a lower bound for $\sum_{N \leq x} R(N)$. If, as (6.58) asserts, $R(N)$ is never too big, then the only way to account for the size of this lower bound is for there to be many terms for which $R(N)$ is nonzero. In other words, $\mathcal{A}$ must be fairly dense. We now make this precise.

Lemma 6.27. As $x \rightarrow \infty$, we have $\sum_{N \leq x} R(N) \gg x^{2} /(\log x)^{2}$.
Proof. By Chebyshev's results from Chapter 3, we have $\pi(x / 2) \gg x / \log x$ as $x \rightarrow \infty$. Thus

$$
\sum_{N \leq x} R(N)=\sum_{N \leq x} \sum_{p+q=N} 1=\sum_{p+q \leq x} 1 \geq\left(\sum_{p \leq x / 2} 1\right)^{2} \gg \frac{x^{2}}{(\log x)^{2}}
$$

Lemma 6.28. As $x \rightarrow \infty$, we have $\sum_{N \leq x} R(N)^{2} \ll x^{3} /(\log x)^{4}$.
Proof. From (6.58),

$$
\begin{aligned}
\sum_{N \leq x} R(N)^{2} & \ll \sum_{2 \leq N \leq x}\left(\frac{N}{(\log N)^{2}} \prod_{p \mid N}\left(1+\frac{1}{p}\right)\right)^{2} \\
& \ll \frac{x^{2}}{(\log x)^{4}} \sum_{2 \leq N \leq x}\left(\prod_{p \mid N}\left(1+\frac{1}{p}\right)\right)^{2} \\
& \ll \frac{x^{2}}{(\log x)^{4}} \sum_{2 \leq N \leq x}\left(\sum_{d \mid N} \frac{1}{d}\right)^{2}
\end{aligned}
$$

It remains to show that the outer sum is $O(x)$. For this, observe that for any natural numbers $d_{1}$ and $d_{2}$,

$$
\left[d_{1}, d_{2}\right] \geq \max \left\{d_{1}, d_{2}\right\} \geq\left(d_{1} d_{2}\right)^{1 / 2}
$$

so that

$$
\begin{aligned}
& \sum_{N \leq x}\left(\sum_{d \mid N} \frac{1}{d}\right)^{2}=\sum_{N \leq x} \sum_{d_{1} \mid N} \sum_{d_{2} \mid N} \frac{1}{d_{1} d_{2}}=\sum_{d_{1}, d_{2} \leq x} \frac{1}{d_{1} d_{2}} \sum_{\substack{N \leq x \\
d_{1}\left|N, d_{2}\right| N}} 1 \\
& \quad \leq \sum_{d_{1}, d_{2} \leq x} \frac{1}{d_{1} d_{2}} \frac{x}{\left[d_{1}, d_{2}\right]} \leq x \sum_{d_{1}, d_{2} \leq x} \frac{1}{\left(d_{1} d_{2}\right)^{\frac{3}{2}}} \leq x\left(\sum_{d=1}^{\infty} d^{-\frac{3}{2}}\right)^{2} \ll x
\end{aligned}
$$

Proof of Theorem 6.26. Writing $R(N)=R(N) \cdot 1$, the Schwarz inequality and Lemmas 6.27 and 6.28 yield that

$$
\begin{aligned}
\frac{x^{4}}{(\log x)^{4}} \ll\left(\sum_{N \leq x} R(N)\right)^{2} & =\left(\sum_{\substack{N \leq x \\
R(N)>0}} R(N) \cdot 1\right)^{2} \\
& \leq \sum_{\substack{N \leq x \\
R(N)>0}} R(N)^{2} \sum_{\substack{N \leq x \\
R(N)>0}} 1 \ll \frac{x^{3}}{(\log x)^{4}} A(x),
\end{aligned}
$$

so that $A(x) \gg x$ as $x \rightarrow \infty$. In other words, $\mathcal{A}$ has positive lower density.

## Notes

The results of this chapter barely begin to scratch the surface of modern sieve theory. Encyclopedic accounts of this subject include the monographs of Halberstam \& Richert [HR74] and Greaves [Gre01]. The introductory texts of Schwarz [Sch74] and Cojocaru \& Murty [CM06] take a more discursive approach. Another treatment of the Brun-Hooley sieve can be found in the the introduction to analytic number theory written by Bateman \& Diamond [BD04].

Excellent references for additive number theory include Ostmann's twovolume work [Ost56] and Nathanson's book [Nat96]. Nathanson's text includes a proof of the following theorem of Vinogradov which should be compared with Theorem 6.22:
$\star$ Theorem 6.29 (Three primes theorem). Let $R_{3}(N)$ denote the number of ways of writing $N$ as an ordered sum of three primes. As $N \rightarrow \infty$ through odd integers, we have

$$
R_{3}(N) \sim \prod_{p}\left(1+\frac{1}{(p-1)^{3}}\right) \prod_{p \mid N}\left(1-\frac{1}{p^{2}-3 p+3}\right) \frac{N^{2}}{2(\log N)^{3}} .
$$

In particular, every sufficiently large odd integer is a sum of three primes.
It follows from Vinogradov's result that every large enough natural number is the sum of at most 4 primes. While Vinogradov's theorem has a similar flavor to Theorem 6.22, the proof, which depends on the circle method, requires substantially deeper input from prime number theory.

See [KT05] for a thorough survey of additive prime number theory.

## Exercises

1. (Gandhi [Gan71], Golomb [Gol74]) For each set of natural numbers $S$, put $w(S):=\sum_{n \in S} 2^{-n}$. For each natural number $k$, let $p_{k}$ denote the $k$ th prime.
(a) If $S$ is the set of natural numbers coprime to $p_{1} \cdots p_{k}$, show that $w(S)=\frac{1}{2}+\frac{1}{2^{p} p^{p+1}}+E$ where $0<E<\frac{1}{2^{p_{k+1}}}$.
(b) Show that for the set $S$ in (a), we have $w(S)=\sum_{d \mid p_{1} \cdots p_{k}} \frac{\mu(d)}{2^{d}-1}$.
(c) Deduce that $p_{k+1}$ is the unique integer for which

$$
1<2^{p_{k+1}}\left(\sum_{d \mid p_{1} \cdots p_{k}} \frac{\mu(d)}{2^{d}-1}-\frac{1}{2}\right)<2
$$

2. (Cf. Nagell [Nag22, §3])
(a) Let $D$ be an integer that is not a square. Using the law of quadratic reciprocity, prove that there is a collection $S$ (say) of $\frac{1}{2} \phi(4|D|)$ residue classes modulo $4|D|$ with the property that for each prime $p \nmid 4 D,\left(\frac{D}{p}\right)=1 \Longleftrightarrow p \bmod 4|D| \in S$.
(b) Deduce from (a) and the results of Chapter 4 that

$$
\sum_{p \leq x,\left(\frac{D}{p}\right)=1} \frac{\log p}{p}=\frac{1}{2} \log x+O(1)
$$

where the implied constant may depend on $D$. (Thus, in a certain average sense, $D$ is a square modulo precisely $\frac{1}{2}$ of all primes.)
(c) Let $F(T)$ be a quadratic polynomial with integer coefficients. Using the sieve of Eratosthenes-Legendre, show that as $x \rightarrow \infty$, the number of $n \leq x$ with $|F(n)|$ prime is $<_{F} x / \log \log x$. (The case when $F(T)=T^{2}+1$ is the third example of $\S 3.2$; cf. Exercise 22.)
3. Use the inclusion-exclusion principle to establish each of the following assertions about squarefree numbers:
(a) The number of squarefree $n \leq x$ is asymptotic to $\frac{1}{\zeta(2)} x=\frac{6}{\pi^{2}} x$ as $x \rightarrow \infty$.
(b) The number of pairs of squarefree integers $n, n+2$ with $1 \leq n \leq x$ is asymptotic to $x \prod_{p}\left(1-2 / p^{2}\right)$ as $x \rightarrow \infty$.
(c) The number of ordered representations of a natural number $N$ as a sum of two positive squarefree integers is asymptotic to

$$
N \prod_{p}\left(1-\frac{2}{p^{2}}\right) \prod_{p^{2} \mid N} \frac{p^{2}-1}{p^{2}-2} \quad(N \rightarrow \infty)
$$

Hint: For each of (a)-(c), first sieve out the multiples of $p^{2}$ for $p \leq z$, where $z=z(x) \rightarrow \infty$ slowly enough to keep the error term in check. To conclude, observe that almost no $n$ are divisible by $p^{2}$ for some prime $p>z$, since $\sum_{p>z} \frac{1}{p^{2}}$ is $o(1)$.
4. (Rényi [Rén55])
(a) Show that for each fixed integer $j \geq 0$, the set of natural numbers $n$ with $\Omega(n)-\omega(n)=j$ possesses an asymptotic density $d_{j}$ (say). Check that $\sum_{j=0}^{\infty} d_{j}=1$.
(b) Show that for all complex numbers $z$ with $|z|<2$, we have

$$
\sum_{j=0}^{\infty} d_{j} z^{j}=\frac{1}{\zeta(2)} \prod_{p}\left(1-\frac{z}{p+1}\right)\left(1-\frac{z}{p}\right)^{-1} .
$$

5. (Hooley [Hoo76], Rieger [Rie77]) If $m$ is an odd natural number, write $l(m)$ for the order of 2 modulo $m$.
(a) Suppose $m \in \mathbf{N}$ is odd and squarefree and put $M:=\operatorname{lcm}[m, l(m)]$. Show that $n \cdot 2^{n}$ runs through every residue class modulo $m$ exactly $M / m$ times as $n$ runs over the integers $1,2,3, \ldots, M$.
(b) Using the result of (a) and the sieve of Eratosthenes-Legendre, show that the set of $n \in \mathbf{N}$ for which $n \cdot 2^{n}+1$ is prime has density zero. (Primes of the form $n \cdot 2^{n}+1$ are called Cullen primes; the first several examples correspond to $n=1,141,4713,5795,6611$, 18496, 32292.)
6. Let $A$ and $B$ be subsets of the natural numbers defined by

$$
\begin{aligned}
& A=\left\{n: n \mid 2^{k}-1 \text { for some positive integer } k\right\}, \\
& B=\left\{n: n \mid 2^{k}+1 \text { for some positive integer } k\right\} .
\end{aligned}
$$

Prove that $A$ has asymptotic density $\frac{1}{2}$ and $B$ has asymptotic density 0 .
7. (Cf. Luca [Luc06, Problem 190]) Let $F_{n}$ denote the $n$th Fibonacci number, so that $F_{0}=0, F_{1}=1$, and for $n>1, F_{n}=F_{n-1}+F_{n-2}$. Show that the set of $n$ for which $F_{n}$ can be written as a sum of two coprime squares has asymptotic density $1 / 2$.
8. Show that for each $d \in \mathbf{N}$, the set of natural numbers $n$ for which $d \mid \varphi(n)$ has asymptotic density 1 . Deduce that the set of $n$ for which $\operatorname{gcd}(n, \varphi(n))=1$ has density zero.
9. (Continuation; cf. Pillai [Pil29]) Let $\mathcal{V}:=\{\varphi(m): m \in \mathbf{N}\}$ be the image of the Euler $\varphi$-function, and let $V(x)$ be the number of $n \leq x$ belonging to $\mathcal{V}$. Show that $V(x)=o(x)$. Hint: Divide the elements $n$ of $\mathcal{V}$ into two classes, depending on whether or not $n$ has a preimage $m$ with only a "small" number of distinct odd prime divisors.

Remark. Maier \& Pomerance [MP88] showed in 1988 that

$$
V(x)=\frac{x}{\log x} \exp \left((C+o(1))(\log \log \log x)^{2}\right)
$$

for a constant $C=0.81781464640 \ldots$. This improved upon earlier results of Erdős, Hall, and Pomerance. The (somewhat complicated) exact order of magnitude of $V(x)$ was subsequently determined by Ford [For98a, For98b].
10. (Blecksmith, Erdős \& Selfridge [BES99]) Say that a prime $p$ is a cluster prime if every even natural number $n<p-2$ can be written in the form $q-q^{\prime}$, where $q$ and $q^{\prime}$ are primes $\leq p$.
(a) Check (perhaps with the aid of a computer) that every prime $p<97$ is a cluster prime, but that $p=97$ is not.
(b) Show that if $p$ is a cluster prime, then for every integer $3 \leq t \leq p-3$, the number of primes in the closed interval $[p-t, p]$ is $\gg \log t$, where the implied constant is absolute. In other words, the primes to the left of $p$ have to "cluster" around $p$.
(c) Show that contrary to what one might expect from (a), the cluster primes are comparatively rare: For every $k$, the number of cluster primes up to $x$ is $O_{k}\left(x /(\log x)^{k}\right)$ as $x \rightarrow \infty$.
11. (Cf. Erdős $[\mathbf{E r d} \mathbf{3 6}]$ ) For each $r \in \mathbf{N}$, define a function $p_{r}: \mathbf{N} \rightarrow$ \{primes $\} \cup\{\infty\}$ by setting $p_{r}(n)$ equal to the $r$ th smallest prime factor of $n$ if $n$ has at least $r$ distinct prime factors and putting $p_{r}(n)=\infty$ otherwise. Observe that $p_{1}(n)<p_{1}(n+1)$ precisely when $n$ is even. In particular, $p_{1}(n)<p_{1}(n+1)$ on a set of asymptotic density $1 / 2$. Show that for each fixed $r$, we have $p_{r}(n)<p_{r}(n+1)$ on a set of asymptotic density $1 / 2$.

Remark. For each $n>1$, put $P(n)$ equal to the largest prime factor of $n$, and put $P(1)=0$. In the 1930s, Erdős conjectured that $P(n)<$ $P(n+1)$ on a set of asymptotic density $1 / 2$. This remains open. Erdős \& Pomerance have shown that each of the inequalities $P(n)>P(n+1)$ and $P(n)<P(n+1)$ holds for a positive proportion of the natural numbers [EP78].
12. For each prime $p$, let $p^{\prime}$ be the prime immediately following $p$. Show that for each $\epsilon>0$, there is a $K>0$ for which the following holds: For large $x$, all but at most $\epsilon x / \log x$ primes $p \leq x$ satisfy

$$
\frac{1}{K} \log x \leq p^{\prime}-p \leq K \log x .
$$

Remark. It is conjectured (see, e.g., [Sou07, Conjecture 1]) that for each fixed $K>0$, the number of $p \leq x$ with $p^{\prime}-p \leq K \log x$ is asymptotically $\left(1-e^{-K}\right) x / \log x$ as $x \rightarrow \infty$.
13. Call a prime $p M$-reclusive if $|q-p|>M$ for every prime $q \neq p$. Show that for every $M>0$ and every $k \in \mathbf{N}$, there are infinitely many $k$-tuples of consecutive primes all of which are $M$-reclusive. (This strengthens the result of Exercise 4.12.)
14. (Erdős \& Nathanson [EN96]) Let $p_{n}$ be the $n$th prime number (in the usual, increasing order). Use Theorem 6.19 to show that for each $\lambda>2$, the series

$$
\sum_{n=1}^{\infty} \frac{1}{n(\log \log 3 n)^{\lambda}\left(p_{n+1}-p_{n}\right)}
$$

converges. It is conjectured that this result is the best possible, in the sense that the series diverges when $\lambda=2$.
15. For each even natural number $N$, let $R^{*}(N)$ be the number of unordered representations of $N$ as a sum of two primes. Then

$$
R^{*}(N) \leq \pi(N-2)-\pi((N-1) / 2)
$$

with equality holding exactly when $N-p$ is prime for each prime $p$ with $N / 2 \leq p \leq N-2$. Use the estimate (3.21) in conjunction with Theorem 6.17 to prove that this upper bound is attained for only finitely many $N$.

Remark. It has been shown by Deshouillers et al. [DGNP93] that $N=210$ is the largest value for which the upper bound is achieved.
16. By modifying the argument of $\S 5.5$, show that the number of representations of an even natural number $N$ as a sum of two 7 -almost primes is $\gg \frac{N}{(\log N)^{2}} \prod_{p \mid N, p>2} \frac{p-1}{p-2}$, as $N \rightarrow \infty$.
17. (Brun) Prove the following theorems of Brun, announced in [Bru19b]:
(a) Every infinite arithmetic progression $a \bmod m$ with $\operatorname{gcd}(a, m)=1$ contains infinitely many 5 -almost primes. (Naturally, Dirichlet's theorem is off-limits here.)
(b) If $x$ is sufficiently large, there is always an 11-almost prime in the interval $(x, x+\sqrt{x}$.
Suggestion: Imitate the lower bound applications of the text, including the selection of the first several $m_{j}$ by the greedy algorithm, but begin instead with the values $K=2.49, K_{1}=2.50$.
18. (A general version of Brun's method) Fix a natural number $k$.
(a) Let $A>0$. Suppose that to each prime $p \leq x^{A}$, we associate $k_{p} \leq k$ residue classes modulo $p$. Show that the number of natural numbers $n \leq x$ avoiding all of these residue classes is

$$
<_{k, A} x \prod_{p \leq x^{A}}\left(1-\frac{k_{p}}{p}\right) \quad(\text { for } x>0),
$$

where the implied constant is independent of the particular choice of residue classes.
(b) Show that there is a constant $B>0$, depending only on $k$, with the following property: If we choose $k_{p} \leq k$ residue classes modulo $p$ for each prime $p \leq x^{B}$, then the number of natural numbers $n \leq x$ avoiding all these classes is

$$
>_{k} x \prod_{p \leq x^{B}}\left(1-\frac{k_{p}}{p}\right) \quad(\text { for } x \rightarrow \infty),
$$

again uniformly in the particular choice of residue classes.
Hint: Use the Chinese remainder theorem to construct a polynomial $F$ for which $p \mid F(n)$ precisely when $n$ falls into one of the $k_{p}$ chosen residue classes $\bmod p$.

Remark. From (a) and (b) we may rederive the results given in the text regarding the twin prime and Goldbach problems, with a slight loss of precision (in that in our lower bound applications, we obtain $r$-almost primes with an unspecified constant $r$ in place of $r=7$ ). For the twin prime problem, the forbidden residue classes are 0 and $-2 \bmod p$. For the Goldbach problem, the forbidden classes are 0 and $N \bmod p$.

When one sees references to "Brun's method" in the literature, often the author has the results of (a) and (b) in mind.
N. B. The results of Problem 18 suffice to handle all the sieving situations that arise in the remaining exercises in this chapter.
19. Suppose that $y=y(x)$ is a positive-valued function of $x$ for which $\frac{\log y}{\log x} \rightarrow$ 0 as $x \rightarrow \infty$. Show that as $x \rightarrow \infty$, all but $o(x)$ of the natural numbers $n \leq x$ have a prime factor $>y$. In other words, $\Psi(x, y)=o(x)$.
20. (Hardy \& Littlewood [HL23]) Show that $\pi(y+x)-\pi(y) \ll \frac{x}{\log x}$ for $y \geq 0$ and $x \geq 2$, where the implied constant is absolute.
21. ("Brun-Titchmarsh inequality" [Tit30]) Let $x \geq 2$. Suppose that $a$ and $m$ are coprime integers with $1 \leq m<x$. Prove that

$$
\pi(x ; m, a) \ll \frac{x}{\varphi(m) \log \frac{x}{m}},
$$

where the implied constant is absolute. (Recall that $\pi(x ; m, a)$ denotes the number of primes $p \leq x$ with $p \equiv a(\bmod m)$.) Is this still true without the assumption that $a$ and $m$ are relatively prime?
22. Suppose $F(T) \in \mathbf{Z}[T]$ is irreducible over $\mathbf{Q}$ and that the leading coefficient of $F(T)$ is positive. For each natural number $d$, let $\nu(d)$ denote the number of roots of $F$ modulo $d$.
(a) A theorem of Landau (cf. [Lan02, eq. (67)]) asserts that for $x \geq 3$,

$$
\sum_{p \leq x} \frac{\nu(p)}{p}=\log \log x+C_{F}+O_{F}\left(\frac{1}{\log x}\right)
$$

where $C_{F}$ is a constant depending on $F$. Deduce from this result and the Brun-Hooley sieve that the number of $n \leq x$ for which $F(n)$ is prime is $<_{F} x / \log x$, again for $x \geq 3$.
(b) Now impose the additional hypothesis that there is no prime $p$ that divides $F(n)$ for every $n \in \mathbf{Z}$. Show that there is an $r \in \mathbf{N}$, depending only on the degree $g$ of $F$, with the property that $F(n)$ is an $r$-almost prime for infinitely many natural numbers $n$.

Remark. Richert [Ric69] has shown that one can take $r=g+1$.
23. (Yang [Yan82]; see also Webb [Web70]) Using the identities

$$
\frac{4}{n}= \begin{cases}\frac{1}{n(k+1) k}+\frac{1}{n(k+1)}+\frac{1}{q k} & \text { if } n=(4 k-1) q, \\ \frac{1}{n k}+\frac{1}{n q k}+\frac{1}{q k} & \text { if } n+1=(4 k-1) q, \\ \frac{1}{n k}+\frac{1}{n k(q k-1)}+\frac{1}{q k-1} & \text { if } n+4=(4 k-1) q, \\ \frac{1}{n k}+\frac{1}{k(q k-n)}+\frac{1}{n(q k-n)} & \text { if } 4 n+1=(4 k-1) q,\end{cases}
$$

show that the number of $n \leq x$ for which (6.22) is unsolvable is $\ll$ $x /(\log x)^{2}$ as $x \rightarrow \infty$. Deduce that the sum of the reciprocals of all $n$ of this kind converges.
Remark. Vaughan [Vau70] has shown that the number of $n \leq x$ for which (6.22) is unsolvable is $\ll x \exp \left(-c(\log x)^{2 / 3}\right)$ for a positive constant $c$.
24. (Erdős [Erd35c]) In Exercise 3.23, we proved that a typical natural number $n \leq x$ has about $\log \log x$ prime factors. One may wonder whether such a result continues to hold if one restricts $n$ to certain special classes of numbers. Here we treat numbers of the form $p-1$, where $p$ is prime. (Such numbers are important, for example, in the study of the Euler $\varphi$-function.) We show that we do indeed have such a result, and that in fact for each $\epsilon>0$,

$$
\#\{p \leq x:|\omega(p-1)-\log \log x|>\epsilon \log \log x\}<_{\epsilon} x /(\log x)^{1+\delta}
$$

where $\delta>0$ depends on $\epsilon$.
(a) Assume $x \geq 3$. Show that all but $O\left(x /(\log x)^{2}\right)$ natural numbers $n \leq x$ possess both of the following properties:
(i) the largest prime factor $P(n)$ (say) of $n$ satisfies $P(n)>$ $x^{1 /(6 \log \log x)}$,
(ii) $n$ is not divisible by $P(n)^{2}$,

Hint: Use the result of Exercise 3.32 to handle condition (i).
(b) For each nonnegative integer $k$, let $N_{k}$ be the number of primes $p \leq x$ for which $p-1$ has both properties (i) and (ii) and satisfies $\omega(p-1)=k$. Show that

$$
N_{k} \leq \sum_{\substack{a \leq x^{1-1 /(6 \log \log x)} \\ \omega(a)=k-1}} \sum_{\substack{p \leq x \\ a \mid p-1 \text { and } \frac{p-1}{a} \text { is prime }}} 1 .
$$

(c) Show that for each natural number $a<x$,

$$
\sum_{\substack{p \leq x \\ a \mid p-1 \text { and } \frac{p-1}{a} \text { is prime }}} 1 \ll \frac{x}{\varphi(a)\left(\log \frac{x}{a}\right)^{2}},
$$

with an absolute implied constant.
(d) Convince yourself that

$$
\sum_{\substack{a \leq x \\ \omega(a)=k-1}} \frac{1}{\varphi(a)} \leq \frac{1}{(k-1)!}\left(\sum_{p^{l} \leq x} \frac{1}{\varphi\left(p^{l}\right)}\right)^{k-1}
$$

where the right-hand sum is over primes and prime powers $p^{l} \leq x$.
(e) Show that for a certain absolute constant $C$,

$$
N_{k} \ll \frac{x(\log \log x)^{2}}{(\log x)^{2}} \frac{(\log \log x+C)^{k-1}}{(k-1)!},
$$

uniformly in $k$. Complete the proof by summing this estimate over $k<(1-\epsilon) \log \log x$ and $k>(1+\epsilon) \log \log x$.
25. (Erdős, ibid.) Prove that if $\epsilon>0$ is sufficiently small, then the following holds: As $x \rightarrow \infty$, there are $\gg_{\epsilon} x / \log x$ primes $p \leq x$ for which the largest prime divisor of $p+1$ is bounded by $x^{1-\epsilon}$. Hint: If $p+1 \leq x$ and $p+1$ has a prime divisor at least $x^{1-\epsilon}$, then $(p+1) / a$ is prime for some natural number $a \leq x^{\epsilon}$.
26. (Luca $[\mathbf{L u c} \mathbf{0 7}])$ Show that the number of natural numbers not exceeding $x$ which can be written in the form $p^{2}-q^{2}$, where $p$ and $q$ are primes, is $\ll x / \log x$.
27. Call the natural number $n$ twinnish if $d+n / d+1$ is prime for every $d$ dividing $n$. If $p$ is the smaller member of a twin prime pair, then $p$ is twinnish, but there are many other such $n$, for example $n=21$ and (less obviously) $n=190757=7^{2} \cdot 17 \cdot 229$. Prove or disprove: $\sum \frac{1}{n}<\infty$, where the sum is extended over all twinnish numbers $n$.
28. (Hardy \& Littlewood [HL23], cf. Landau [Lan00]) Let $R(N)$ be the number of ordered representations of $N$ as a sum of two primes. Conjecture 3.19 asserts that as $N \rightarrow \infty$ through even numbers,

$$
\begin{equation*}
R(N)=(A+o(1))\left(\prod_{p \mid N} \frac{p-1}{p-2}\right) \frac{N}{(\log N)^{2}}, \tag{6.59}
\end{equation*}
$$

where

$$
\begin{equation*}
A=2 \prod_{p>2}\left(1-\frac{1}{(p-1)^{2}}\right) \tag{6.60}
\end{equation*}
$$

This differs from what a naive sieve argument would suggest, namely that (6.59) holds with

$$
\begin{equation*}
A=8 \exp (-2 \gamma) \prod_{p>2}\left(1-\frac{1}{(p-1)^{2}}\right) \tag{6.61}
\end{equation*}
$$

In this exercise we outline a proof that (6.61) cannot be correct. In fact, we show that if an asymptotic relation of the form (6.59) holds, then $A$ must be given by (6.60).
(a) Use the prime number theorem to show that $\sum_{N \leq x} R(N) \sim \frac{1}{2} \frac{x^{2}}{(\log x)^{2}}$ as $x \rightarrow \infty$.
(b) Deduce from (a) that as $x \rightarrow \infty$,

$$
\sum_{2 \leq N \leq x} \frac{R(N)}{N /(\log N)^{2}} \sim x .
$$

(c) Put $g(N):=\prod_{p \mid N, p>2} \frac{p-1}{p-2}$ for each $N$, and define an arithmetic function $h$ by the relation $g(N)=\sum_{d \mid N} h(d)$. Show that $h$ is supported on odd, squarefree positive integers, and that as $x \rightarrow \infty$,

$$
\frac{1}{x} \sum_{\substack{N \leq x \\ N \text { even }}} g(N) \rightarrow \frac{1}{2} \sum_{d \text { odd }} \frac{h(d)}{d}=\frac{1}{2} \prod_{p>2} \frac{(p-1)^{2}}{p(p-2)} .
$$

(d) Use the result of (c) and the purported relation (6.59) to derive another asymptotic formula for $\sum_{2 \leq N \leq x} \frac{R(N)}{N /(\log N)^{2}}$ which, when compared with that of (b), proves (6.60).

Remark. The methods used to prove Vinogradov's three primes theorem can be employed to show that in fact the relation (6.59) with $A$ given by (6.60) holds for almost all even natural numbers $N$ (see, e.g, [Vau97, §3.2]). More precisely, (6.59) holds (with this $A$ ) as $N \rightarrow \infty$ through even numbers, provided we exclude a particular set of even numbers $N$ of asymptotic density zero.
29. (Landau [Lan30]) Show that under the hypotheses of Theorem 6.23, the set $\mathcal{A}$ is a basis of order at most $2\lfloor 1 / \delta(\mathcal{A})\rfloor$.
30. Say that a set $\mathcal{A} \subset \mathbf{N}_{0}$ is an asymptotic basis of finite order if $\mathbf{N} \backslash h \mathcal{A}$ is finite for some $h \in \mathbf{N}$.
(a) Show that if $a_{1}, \ldots, a_{k} \in \mathbf{N}$ and $\operatorname{gcd}\left(a_{1}, \ldots, a_{k}\right)=1$, then every sufficiently large natural number can be written in the form $\sum_{i=1}^{k} a_{i} x_{i}$, where each $x_{i} \in \mathbf{N}_{0}$.
(b) Let $\mathcal{A}$ be a subset of $\mathbf{N}_{0}$. Suppose that $0 \in \mathcal{A}$, that $\mathcal{A}$ has positive lower density (i.e., (6.57) holds), and that there is no integer $d>1$ dividing each $a \in \mathcal{A}$. Show that $\mathcal{A}$ is an asymptotic basis of finite order.
31. (Landau, ibid.; see also Nathanson [Nat87a]) Suppose $\mathcal{P}$ is a set of primes with the property that

$$
\liminf _{x \rightarrow \infty} \frac{\#\{p \in \mathcal{P}: p \leq x\}}{x / \log x}>0
$$

Show that there is a constant $S_{\mathcal{P}}$ with the property that every sufficiently large natural number is the sum of at most $S_{\mathcal{P}}$ primes all of which belong to $\mathcal{P}$.
32. (Prachar [Pra52]) Show that for large $x$, there are $\gg x$ natural numbers $n \leq x$ that can be written in the form $q-p$, where $p, q \leq x$ are primes. Hint: Adapt the second-moment method appearing in the proof of Schnirelmann's theorem.
33. (Continuation) For each prime $p$, write $p^{\prime}$ for the prime immediately following $p$. Show that for some constant $K>0$, the following holds: For all large $x$, there are $\gg \log x$ natural numbers $n \leq K \log x$ which can be written in the form $p^{\prime}-p$ for some prime $p \leq x$. Hint: Use Exercise 12.
34. (Romanov [Rom34]) Let $r(n)$ be the number of representations of $n$ in the form $2^{k}+p$, where $p$ is prime and $k \geq 1$. In this exercise and the next, we sketch a proof that $r(n)>0$ on a set of positive lower density. In Exercise 36, we prove the complementary result that $r(n)=0$ on a set of odd numbers of positive density.
(a) Show that for all natural numbers $n$, we have $\sum_{d \mid n} \frac{1}{d} \ll \log \log 3 n$.
(b) For each odd integer $d$, let $l(d)$ denote the order of 2 modulo $d$. Show that if $l(d) \leq x$, then $d$ divides $D:=\prod_{1 \leq k \leq x}\left(2^{k}-1\right)$. Deduce from (a) that $\sum_{l(d) \leq x} d^{-1} \ll \log (2 x)$ for $x \geq 1$.
(c) Using partial summation, prove that $\sum_{\substack{d \text { odd } \\ d \geq 1}} \frac{1}{d \cdot l(d)}<\infty$.
35. (Continuation)
(a) Show that $\sum_{n \leq x} r(n) \gg x$ as $x \rightarrow \infty$.
(b) Show that $\sum_{n \leq x} r(n)^{2}$ does not exceed the number of solutions $\left(p_{1}, p_{2}, k_{1}, k_{2}\right)$ to

$$
p_{2}-p_{1}=2^{k_{1}}-2^{k_{2}}
$$

where $p_{1}, p_{2}$ are primes $\leq x$ and $1 \leq k_{1}, k_{2} \leq \log x / \log 2$.
(c) Show that the number of solutions as in (b) is $\ll x$. Hint: To estimate the number of solutions with $k_{1} \neq k_{2}$, use Theorem 6.19 and the result of Exercise 34(c).
(d) Deduce from (a)-(c), and the Cauchy-Schwarz inequality that there are $\gg x$ natural numbers $n \leq x$ for which $r(n)>0$.
36. (Continuation; Erdős [Erd50b], following [Sie88, Chapter XII])
(a) Check that every integer $k$ belongs to at least one of the congruence classes $0 \bmod 2,0 \bmod 3,1 \bmod 4,3 \bmod 8,7 \bmod 12,23 \bmod 24$.
(b) Suppose $n \equiv 1(\bmod 3), n \equiv 1(\bmod 7), n \equiv 2(\bmod 5), n \equiv$ $2^{3}(\bmod 17), n \equiv 2^{7}(\bmod 13)$, and $n \equiv 2^{23}(\bmod 241)$. Show that for every integer $k \geq 0$, the number $n-2^{k}$ is divisible by some prime from the set $\{3,5,7,13,17,241\}$.
(c) Suppose that in addition to the congruences in (b), we require also that $n \equiv 1(\bmod 2)$ and $n \equiv 3(\bmod 31)$. Show that the positive $n$ satisfying all of these congruences comprise an infinite arithmetic progression of odd integers $n$ with $r(n)=0$.

## An Elementary Proof of the Prime Number Theorem

No elementary proof of the prime number theorem is known, and one may ask whether it is reasonable to expect one. Now we know that the theorem is roughly equivalent to a theorem about an analytic function, the theorem that Riemann's zeta function has no roots on a certain line. A proof of such a theorem, not fundamentally dependent on the theory of functions, seems to me extraordinarily unlikely. It is rash to assert that a mathematical theorem cannot be proved in a particular way; but one thing seems quite clear. We have certain views about the logic of the theory; we think that some theorems, as we say, "lie deep" and others nearer to the surface. If anyone produces an elementary proof of the prime number theorem, he will show that these views are wrong, that the subject does not hang together in the way we have supposed, and that it is time for the books to be cast aside and for the theory to be rewritten. - G. H. Hardy [Boh52]

## 1. Introduction

Recall that the prime number theorem asserts that as $x \rightarrow \infty$,

$$
\begin{equation*}
\pi(x)=(1+o(1)) \frac{x}{\log x} . \tag{7.1}
\end{equation*}
$$

In Chapter 3, we described the early history of this result, including its origin as a conjecture by a young Gauss and its eventual proof by Hadamard and de la Vallée-Poussin (independently) in 1896, following a plan laid out by Riemann. Their proofs relied heavily on results from the then-budding field of complex analysis.

In 1931, Wiener and Ikehara proved the following theorem, which leads quickly to a proof of the prime number theorem requiring only scant knowledge of the analytic properties of the Riemann zeta-function $\zeta(s)$ :
$\star$ Theorem 7.1. Let $\sum_{n=1}^{\infty} f(n) n^{-s}$ be a Dirichlet series with nonnegative coefficients, convergent for $\Re(s)>1$. Let $F$ be the (analytic) function defined by the series in this region, and suppose that $F$ can be extended to a function analytic on an open set containing $\Re(s) \geq 1$, except possibly for a simple pole at $s=1$. If $R$ is the residue of $F$ at $s=1$, then

$$
\sum_{n \leq x} f(n)=(R+o(1)) x \quad(x \rightarrow \infty) .
$$

Let us briefly sketch the derivation of the prime number theorem from Theorem 7.1. An easy calculation (Exercise 1) shows that

$$
\begin{equation*}
\zeta(s)=1+\frac{1}{s-1}-s \int_{1}^{\infty} \frac{\{x\}}{x^{s+1}} d x \tag{7.2}
\end{equation*}
$$

in the region $\Re(s)>1$. The integral in (7.2) is analytic for $\Re(s)>0$, and so $\zeta(s)$ can be continued to a function which is analytic for $\Re(s)>0$, except for a simple pole at $s=1$ with residue 1 . Since $\zeta(s)$ has no zeros for $\Re(s)>1$ (since it can be written as an absolutely convergent Euler product there), if one can show that $\zeta(s)$ also has no zeros on $\Re(s)=1$, then $-\zeta^{\prime}(s) / \zeta(s)$ analytically continues to an open set containing $\Re(s) \geq 1$, apart from a simple pole at $s=1$ with residue 1 . Since

$$
\begin{array}{r}
\zeta(s)=\prod_{p}\left(1-\frac{1}{p^{s}}\right)^{-1}, \quad \text { we obtain by logarithmic differentiation }  \tag{7.3}\\
-\frac{\zeta^{\prime}(s)}{\zeta(s)}=\sum_{p} \frac{\log p / p^{s}}{1-1 / p^{s}}=\sum_{p}\left(\frac{\log p}{p^{s}}+\frac{\log p}{p^{2 s}}+\cdots\right)=\sum_{n=1}^{\infty} \frac{\Lambda(n)}{n^{s}},
\end{array}
$$

and so the Wiener-Ikehara result shows that

$$
\psi(x)=\sum_{n \leq x} \Lambda(n)=(1+o(1)) x
$$

an assertion we have seen to be equivalent to the prime number theorem (Corollary 3.8). Conversely, if $\zeta(s)$ has any zeros on the line $\Re(s)=1$, then it is relatively easy to prove directly that the prime number theorem cannot hold (see Exercise 4).

Thus the prime number theorem is, more or less, equivalent to an analytic assertion, namely the nonvanishing of $\zeta(s)$ on the line $\Re(s)=1$. How could an elementary, real-variables proof establish an inherently complexanalytic fact such as this? It was this line of reasoning that led many prominent mathematicians, including Hardy, to the mistaken conclusion that such an elementary proof probably did not exist. When such a proof surfaced in 1948, it sent shockwaves throughout the world of mathematics.
1.1. Selberg's fundamental formula and its consequences. The key ingredient in the early elementary proofs of the prime number theorem is the fundamental formula (also called the symmetry formula) discovered by Selberg in March of 1948,

$$
\begin{equation*}
\theta(x) \log x+\sum_{p \leq x} \theta\left(\frac{x}{p}\right) \log p=2 x \log x+O(x) . \tag{7.4}
\end{equation*}
$$

The proof, which appears below in $\S 3$, can be understood by a talented high-school student. But the implications of this formula are unexpectedly far-reaching. One striking consequence was noticed by Selberg early on (already by April of 1948). Chebyshev had shown (see Theorem 3.5) that

$$
0<a:=\lim \inf \frac{\theta(x)}{x} \leq \lim \sup \frac{\theta(x)}{x}=: A<\infty .
$$

Using the symmetry formula, one can effect a simple proof that

$$
A+a=2,
$$

a result not easily accessible to other elementary methods. Indeed, let $x \rightarrow$ $\infty$ along a sequence of values on which $\theta(x)=(A+o(1)) x$. Then for the left-hand side of (7.4) we have the estimate

$$
\begin{aligned}
\theta(x) \log x+\sum_{p \leq x} \theta & \left(\frac{x}{p}\right) \log p \geq(A+o(1)) x \log x \\
& +\sum_{p \leq x / \log x}\left((a+o(1)) \frac{x}{p}\right) \log p=(A+a+o(1)) x \log x,
\end{aligned}
$$

so that (7.4) implies $A+a \leq 2$. If we begin instead with a sequence on which $\theta(x)=(a+o(1)) x$, then a similar argument yields the reverse inequality $A+a \geq 2$. So $A+a=2$.

In July, 1948, Turán gave a seminar at the Institute for Advanced Study on Selberg's elementary proof of Dirichlet's theorem on primes in progressions. In passing, he mentioned Selberg's fundamental formula. Erdős, who was in the audience, quickly realized that (7.4) could be used to give an elementary proof that the ratio $p_{n+1} / p_{n}$ of consecutive primes tends to 1 . Actually Erdős was able to deduce from Selberg's formula the stronger result that for any $\delta>0$, there are $>c(\delta) x / \log x$ primes in the interval $(x,(1+\delta) x]$ (for sufficiently large $x$ ).

Erdős excitedly described his result and proof to Selberg. Two days later, on July 18, 1948, Selberg used Erdős's result to fashion the first elementary proof of the prime number theorem. Selberg's original argument and certain simplifications, due to Selberg and Erdős, are described in [Erd49].
1.2. Proving the prime number theorem from the symmetry formula. The proof of the prime number theorem given in this chapter is similar to the one ultimately published by Selberg in the Annals of Mathematics [Sel49b]. Define the remainder term $R(x)$ by the formula $\theta(x)=x+R(x)$, so that the prime number theorem is equivalent to the estimate $R(x)=o(x)$. From the fundamental formula (7.4) one easily deduces (cf. (7.27)) that

$$
\begin{equation*}
|R(x)| \log x \leq \sum_{p \leq x}|R(x / p)| \log p+O(x) . \tag{7.5}
\end{equation*}
$$

The prime number theorem says that $\sum_{p \leq x} \log p \sim x$, so that (7.5) should translate under partial summation to an estimate of the shape

$$
|R(x)| \log x \lesssim \sum_{n \leq x}|R(x / n)| .
$$

It turns out that an estimate of this kind can be deduced starting from the fundamental formula without appeal to the prime number theorem, namely

$$
\begin{equation*}
|R(x)| \log x \leq \sum_{n \leq x}|R(x / n)|+O(x \log \log 3 x) . \tag{7.6}
\end{equation*}
$$

(See (7.34).) This is more convenient to work with than (7.5), because in (7.6) the primes do not explicitly appear on the right-hand side.

Let us suppose that $\alpha:=\lim \sup _{x \rightarrow \infty}|R(x)| / x$. Then $\alpha<\infty$, since $\theta(x) \ll x$, and the prime number theorem is the assertion that $\alpha=0$. From (7.6), we find that

$$
\frac{|R(x)|}{x} \lesssim \frac{1}{x \log x} \sum_{n \leq x}|R(x / n)| \lesssim \frac{1}{x \log x} \sum_{n \leq x} \alpha \frac{x}{n} \approx \alpha .
$$

In fact, if one is a little careful here, one gets from this argument that

$$
\begin{equation*}
\lim \sup |R(x)| / x \leq \alpha \tag{7.7}
\end{equation*}
$$

Given how we defined $\alpha$, the reader will be forgiven if she is not impressed by (7.7)! But there is reason to take heart: Granted, (7.7) doesn't tell us anything that we don't already know; in the words of H. N. Shapiro [Sha83], (7.6) is a balanced inequality, meaning that it returns whatever upper bound on $\lim \sup |R(x)| / x$ that it is fed. But if the right-hand side of (7.7) had been any smaller, we would have a contradiction to the choice of $\alpha$. The plan of the proof is to show that unless $\alpha=0$, one can indeed get an upper bound for $\lim \sup |R(x)| / x$ improving upon $\alpha$. This contradiction forces us to have $\alpha=0$, so that the prime number theorem follows.

Actually the means of producing such an improvement are a bit clearer if we part ways from Selberg and work with integrals instead of sums. (This approach seems to have been introduced by Wright [Wri52]. The similar approach we take here is due to Nevanlinna [Nev62].) Rescale the remainder term $R(x)$ by introducing the function $r(x):=e^{-x} R\left(e^{x}\right)$. Then the prime number theorem amounts to the assertion that $r(x)=o(1)$. Instead of working with (7.6), we work with the corresponding integral inequality

$$
\begin{equation*}
|r(x)| \leq \int_{0}^{x}|r(t)| d t+o(1) . \tag{7.8}
\end{equation*}
$$

(See Theorem 7.10.) In parallel with the above, if we suppose lim sup $|r(x)|=$ $\lambda$, then (7.8) returns to us to the same estimate. In order to forcibly unbalance the inequality (7.8), Nevanlinna examines what happens between the sign changes of $r(x)$, showing that if $\lambda>0$, then over each interval between sign changes, $|r(x)|$ is quite often appreciably smaller than $\lambda$. This implies that (7.8) returns an improved estimate unless $\lambda=0$. Thus $r(x)=o(1)$.

Notation. If $A$ is a bounded subset of $\mathbf{R}$, the expression $\int_{A} f(t) d t$ should be read as a synonym for the (improper) Riemann integral $\int_{-\infty}^{\infty} \chi_{A}(t) f(t) d t$, where $\chi_{A}$ is the indicator function of $A$. The (Jordan) measure $\mu(A)$ of $A$ is defined by $\mu(A):=\int_{A} 1 d t$. When these expressions exist, their values agree with those from the Lebesgue theory of integration, but this chapter can be read without any knowledge of that subject.

When any of $p, q$, and $r$ appear in the conditions of summation in this chapter, they always denote primes.

## 2. Chebyshev's theorems revisited

Recall the following three results from Chapter 3: First, $\pi(x) \ll x / \log x$. Second, $\pi(x) \gg x / \log x$. Third, if there is a constant $C$ for which $\pi(x)=$ $(C+o(1)) x / \log x$, then necessarily $C=1$. Our approach to the Selberg symmetry formula will be clearer if we first revisit these results of Chebyshev from a somewhat different perspective.

In Chapter 3, the identity $\sum_{d \mid n} \Lambda(d)=\log n$ played the key role. If we Möbius-invert this identity, we find that

$$
\begin{equation*}
\Lambda(n)=\sum_{d \mid n} \mu(d) \log \frac{n}{d}=\sum_{a b=n} \mu(a) \log b \tag{7.9}
\end{equation*}
$$

Thus

$$
\begin{align*}
\psi(x)=\sum_{n \leq x} \Lambda(n) & =\sum_{a b \leq x} \mu(a) \log b \\
& =\sum_{a \leq x} \mu(a)\left(\frac{x}{a} \log \frac{x}{a}-\frac{x}{a}+O\left(\log \frac{e x}{a}\right)\right) \tag{7.10}
\end{align*}
$$

using Lemma 3.10 to estimate $\sum_{b \leq x / a} \log b$. (Here $e$ is the usual base of the natural logarithm. The factor of $e$ is included in the error term so that the estimate is valid even when $x / a$ is very close to 1.) This does not look like a promising approach to estimating $\psi(x)$, because at this point we have no way to estimate the sums of the Möbius function that appear. But as we will see shortly, this barrier is not at all insurmountable.

### 2.1. Another Möbius inversion formula.

Lemma 7.2. Let $f$ and $g$ be any two complex-valued functions on $[1, \infty)$ satisfying the functional equation

$$
f(x)=\sum_{n \leq x} g(x / n)
$$

Then

$$
g(x)=\sum_{n \leq x} \mu(n) f(x / n)
$$

Proof. If $f$ and $g$ obey the given relation, then

$$
\begin{aligned}
\sum_{n \leq x} \mu(n) f(x / n) & =\sum_{n \leq x} \mu(n) \sum_{m \leq x / n} g\left(\frac{x}{m n}\right) \\
& =\sum_{m n \leq x} \mu(n) g\left(\frac{x}{m n}\right)=\sum_{N \leq x} g\left(\frac{x}{N}\right) \sum_{m \mid N} \mu(m)=g(x)
\end{aligned}
$$

since $\sum_{m \mid N} \mu(m)$ vanishes unless $N=1$.
Remark. If $f$ and $g$ are arithmetic functions, we may extend their domain to $[1, \infty)$ by declaring that they vanish at nonintegral arguments. Then Lemma 7.2 reduces to one direction of the usual Möbius inversion formula.
Corollary 7.3. For $x \geq 1$,
(i) $\sum_{n \leq x} \frac{\mu(n)}{n}=O(1)$,
(ii) $\sum_{n \leq x} \frac{\mu(n)}{n} \log \frac{x}{n}=O(1)$,
(iii) $\sum_{n \leq x} \frac{\mu(n)}{n}\left(\log \frac{x}{n}\right)^{2}=2 \log x+O(1)$.

Proof. We apply the inversion formula of Lemma 7.2 for three different choices of $f$ and $g$. First, take $g$ to be identically 1. Then $\sum_{n \leq x} g(x / n)=$ $\lfloor x\rfloor$, and so taking $f(x):=\lfloor x\rfloor$, Lemma 7.2 gives us that

$$
1=\sum_{n \leq x} \mu(n)\lfloor x / n\rfloor=\sum_{n \leq x} \mu(n)\left(\frac{x}{n}+O(1)\right),
$$

from which (i) easily follows. Next, apply Lemma 7.2 with $g(x):=x$ and $f(x):=\sum_{n \leq x} x / n$. Since $f(x)=x \log x+\gamma x+O(1)$, we find that

$$
x=\sum_{n \leq x} \mu(n)\left(\frac{x}{n} \log \frac{x}{n}+\gamma \frac{x}{n}+O(1)\right) .
$$

Rearranging this estimate and using (i) yields (ii). Lastly, take $g(x):=$ $x \log x$ and $f(x):=\sum_{n \leq x} g(x / n)$. Then

$$
\begin{aligned}
f(x) & =\sum_{n \leq x} \frac{x}{n} \log \frac{x}{n} \\
& =x \log x \sum_{n \leq x} \frac{1}{n}-x \sum_{n \leq x} \frac{\log n}{n} .
\end{aligned}
$$

It is easy to show (by imitating the proof of Theorem 3.16) that

$$
\sum_{n \leq x} \frac{\log n}{n}=\frac{1}{2}(\log x)^{2}+c+O\left(\frac{\log e x}{x}\right)
$$

for some positive constant $c$. Thus

$$
\begin{aligned}
f(x) & =x \log x\left(\log x+\gamma+O\left(\frac{1}{x}\right)\right)-x\left(\frac{1}{2}(\log x)^{2}+c+O\left(\frac{\log e x}{x}\right)\right) \\
& =\frac{1}{2} x(\log x)^{2}+\gamma x \log x-c x+O(\log e x) .
\end{aligned}
$$

So from Lemma 7.2 and (i) and (ii), we find that

$$
\begin{aligned}
x \log x & =\sum_{n \leq x} \mu(n) f(x / n) \\
& =\sum_{n \leq x} \mu(n)\left(\frac{1}{2} \frac{x}{n}\left(\log \frac{x}{n}\right)^{2}+\gamma \frac{x}{n} \log \frac{x}{n}-c \frac{x}{n}+O\left(\log \frac{e x}{n}\right)\right) \\
& =\frac{1}{2} x \sum_{n \leq x} \frac{\mu(n)}{n}\left(\log \frac{x}{n}\right)^{2}+O(x)+O\left(\sum_{n \leq x} \log \frac{e x}{n}\right) .
\end{aligned}
$$

The final error term here is also $O(x)$ (cf. (4.23)), so that dividing by $\frac{1}{2} x$ gives us (iii).
2.2. Another proof of Chebyshev's results. With Corollary 7.3 in hand, we can again pick up our new approach to Chebyshev's results. In (7.10), we found that

$$
\psi(x)=x \sum_{a \leq x} \frac{\mu(a)}{a} \log \frac{x}{a}-x \sum_{a \leq x} \frac{\mu(a)}{a}+O\left(\sum_{a \leq x} \log \frac{e x}{a}\right)
$$

and Corollary 7.3 (parts (i) and (ii)) says that both of the sums here are $O(1)$. Since the $O$-term is $O(x)$, this yields another proof that $\psi(x) \ll x$, which is equivalent to the upper estimate $\pi(x) \ll x / \log x$.

What about the latter two results of Chebyshev? Suppose we multiply the identity (7.9) by $1 / n$ before summing; then we obtain

$$
\begin{aligned}
\sum_{n \leq x} \frac{\Lambda(n)}{n} & =\sum_{a b \leq x} \frac{\mu(a)}{a b} \log b=\sum_{a \leq x} \frac{\mu(a)}{a} \sum_{b \leq x / a} \frac{\log b}{b} \\
& =\sum_{a \leq x} \frac{\mu(a)}{a}\left(\frac{1}{2}\left(\log \frac{x}{a}\right)^{2}+c+O\left(\frac{\log e(x / a)}{x / a}\right)\right) \\
& =\frac{1}{2} \sum_{a \leq x} \frac{\mu(a)}{a}\left(\log \frac{x}{a}\right)^{2}+c \sum_{a \leq x} \frac{\mu(a)}{a}+O(1) .
\end{aligned}
$$

Applying (i) and (iii) of Corollary 7.3, we arrive at the estimate

$$
\begin{equation*}
\sum_{n \leq x} \frac{\Lambda(n)}{n}=\log x+O(1) \tag{7.11}
\end{equation*}
$$

The estimate (7.11) by itself can be used to rederive all three of Chebyshev's results. For example, if $\psi(x)=(C+o(1)) x$ for a constant $C$, then partial summation implies that $\sum_{n \leq x} \Lambda(n) / n=(C+o(1)) \log x$, so that we must have $C=1$. Also, from (7.11) we find that one can choose a constant $B$ for
which

$$
\begin{equation*}
\sum_{x<n \leq B x} \frac{\Lambda(n)}{n}>1 \tag{7.12}
\end{equation*}
$$

for all $x \geq 1$. But the left-hand side of (7.12) is bounded above by $\psi(B x) / x$. Hence $\psi(B x)>x$ and so $\psi(x)>x / B$ whenever $x \geq B$. This implies the lower estimate $\pi(x) \gg x / \log x$ as $x \rightarrow \infty$. A similar argument, omitted here, would show that (7.11) by itself also implies the upper estimate $\pi(x) \ll$ $x / \log x$.

The upshot of our work in this section is that the Möbius sum estimates of Corollary 7.3 contain all the information about primes embodied in these three results of Chebyshev. As we shall establish in the remainder of this chapter, the estimates of Corollary 7.3 in fact already contain the prime number theorem.

## 3. Proof of Selberg's fundamental formula

3.1. An identity of arithmetic functions. Our jumping-off point for the proof of Selberg's fundamental formula is the following identity, whose (formal) verification requires only the familiar quotient rule from differential calculus:

$$
\begin{equation*}
\frac{\zeta^{\prime \prime}(s)}{\zeta(s)}=\left(\frac{\zeta^{\prime}(s)}{\zeta(s)}\right)^{\prime}+\left(\frac{\zeta^{\prime}(s)}{\zeta(s)}\right)^{2} . \tag{7.13}
\end{equation*}
$$

To get at the arithmetic content implicit in this identity, we expand both sides of (7.13) as Dirichlet series (in the region $\Re(s)>1$ ) and then equate corresponding coefficients.

This is straightforward once we know how to multiply Dirichlet series. If $f$ is an arithmetic function, let us agree that the Dirichlet series associated to $f$ refers to the function $F$ defined by

$$
F(s):=\sum_{n=1}^{\infty} \frac{f(n)}{n^{s}}
$$

with domain consisting of those complex numbers $s$ for which the series converges. Suppose that $F$ and $G$ are the Dirichlet series associated with $f$ and $g$, respectively, and that the series defining $F$ and $G$ converge absolutely at $s$. Then

$$
\begin{align*}
F(s) G(s) & =\left(\sum_{n=1}^{\infty} \frac{f(n)}{n^{s}}\right)\left(\sum_{m=1}^{\infty} \frac{g(m)}{m^{s}}\right) \\
& =\sum_{n, m \in \mathbf{N}} \frac{f(n) g(m)}{(n m)^{s}}=\sum_{N=1}^{\infty} \frac{h(N)}{N^{s}}, \tag{7.14}
\end{align*}
$$

where

$$
\begin{equation*}
h(N):=\sum_{n m=N} f(n) g(m) \tag{7.15}
\end{equation*}
$$

The function $h$ is referred to as the Dirichlet convolution of $f$ and $g$.
We can now obtain Dirichlet series expansions of both sides of (7.13).
Differentiating $\zeta(s)$ twice, term-by-term, shows that (for $\Re(s)>1$ )

$$
\begin{equation*}
\zeta^{\prime \prime}(s)=\sum_{n=1}^{\infty} \frac{(\log n)^{2}}{n^{s}} \tag{7.16}
\end{equation*}
$$

The Euler product representation of $\zeta(s)$ implies that (for $\Re(s)>1$ )

$$
\frac{1}{\zeta(s)}=\prod_{p}\left(1-\frac{1}{p^{s}}\right)=\sum_{n=1}^{\infty} \frac{\mu(n)}{n^{s}}
$$

So from (7.14) and (7.15), the left-hand side of (7.13) is represented by the Dirichlet series associated to the convolution of $\mu$ and $\log ^{2}$. To handle the right-hand side, we recall from the introduction that for $\Re(s)>1$,

$$
-\frac{\zeta^{\prime}(s)}{\zeta(s)}=\sum_{n=1}^{\infty} \frac{\Lambda(n)}{n^{s}}
$$

From this we easily read off a Dirichlet series expansion of the right-hand side of (7.13); equating this expansion coefficient-by-coefficient with what we obtained for the left-hand side, we find that for each natural number $n$,

$$
\begin{equation*}
\sum_{a b=n} \mu(a)(\log b)^{2}=\Lambda(n) \log n+\sum_{a b=n} \Lambda(a) \Lambda(b) \tag{7.17}
\end{equation*}
$$

This identity of arithmetic functions will be used below, in combination with the results of Corollary 7.3, to prove Selberg's fundamental formula.

But is our derivation of (7.17) legal? By equating coefficients as above, we are implicitly assuming that (under reasonable hypotheses) the same function cannot have two Dirichlet series expansions. We could prove such a result; this is not hard (see, e.g., [Apo76, Theorem 11.3]), but it would take us somewhat afield. Alternatively, it is possible to develop a theory of formal Dirichlet series which allows one to justify all of the above manipulations without any recourse to analysis (see, e.g., [Sha83, Chapter 4]). Again, this would take us somewhat off point. Perhaps the simplest procedure is to view the above argument simply as a heuristic suggesting (7.17). We can then try to prove (7.17) directly.

This last plan is relatively painless to execute. The left-hand side of (7.17) can be rewritten as

$$
\begin{align*}
\sum_{a b=n} \mu(a)\left(\sum_{d \mid b} \Lambda(b)\right)^{2} & =\sum_{d_{1}\left|n, d_{2}\right| n} \Lambda\left(d_{1}\right) \Lambda\left(d_{2}\right) \sum_{\substack{a b=n \\
\left[d_{1}, d_{2}\right] \mid b}} \mu(a) \\
& =\sum_{\substack{d_{1}, d_{2} \\
\left[d_{1}, d_{2}\right] \mid n}} \Lambda\left(d_{1}\right) \Lambda\left(d_{2}\right) \sum_{a \left\lvert\, \frac{n}{\left[d_{1}, d_{2}\right]}\right.} \mu(a) \\
& =\sum_{\substack{d_{1}, d_{2} \\
\left[d_{1}, d_{2}\right]=n}} \Lambda\left(d_{1}\right) \Lambda\left(d_{2}\right), \tag{7.18}
\end{align*}
$$

where $\left[d_{1}, d_{2}\right]$ denotes the least common multiple of $d_{1}$ and $d_{2}$. But the von Mangoldt function $\Lambda$ is supported on prime powers. So to prove (7.17), it is enough to check that (7.18) and the right-hand side of (7.17) agree when $\omega(n)=1$ or 2 , since in all other cases both expressions vanish. But if $n=p^{e}$, then both expressions equal $(2 e-1)(\log p)^{2}$, while if $n=p_{1}^{e_{1}} p_{2}^{e_{2}}$ (with $p_{1} \neq p_{2}$ ), then both come out to $2 \log p_{1} \log p_{2}$.
3.2. Estimating. Starting with the identity (7.17), we sum over $n \leq x$ to find that

$$
\sum_{a b \leq x} \mu(a)(\log b)^{2}=\sum_{n \leq x} \Lambda(n) \log n+\sum_{a b \leq x} \Lambda(a) \Lambda(b) .
$$

We would like an estimate for the left-hand side with an error term of at most $O(x)$. Write

$$
\begin{align*}
& \sum_{a \leq x} \mu(a) \sum_{b \leq x / a}(\log b)^{2}=\sum_{a \leq x} \mu(a)\left(\int_{1}^{x / a}(\log t)^{2} d t+O\left(\left(\log \frac{x}{a}\right)^{2}\right)\right)  \tag{7.19}\\
& \quad=\sum_{a \leq x} \mu(a)\left(\frac{x}{a}\left(\log \frac{x}{a}\right)^{2}-2 \frac{x}{a} \log \frac{x}{a}+2 \frac{x}{a}-2\right)+O\left(\sum_{a \leq x}\left(\log \frac{x}{a}\right)^{2}\right) .
\end{align*}
$$

The error term here is

$$
\ll \int_{1}^{x}\left(\log \frac{x}{t}\right)^{2} d t+O\left((\log x)^{2}\right) \ll x
$$

by a straightforward calculation. The main terms of (7.19) are estimated for us by Corollary 7.3, and collecting these estimates shows that

$$
\begin{equation*}
\sum_{n \leq x} \Lambda(n) \log n+\sum_{a b \leq x} \Lambda(a) \Lambda(b)=2 x \log x+O(x) . \tag{7.20}
\end{equation*}
$$

Now

$$
\begin{align*}
\sum_{n \leq x} \Lambda(n) \log n & =\int_{1}^{x} \log t d \psi(t)  \tag{7.21}\\
& =\psi(x) \log x-\int_{1}^{x} \frac{\psi(t)}{t} d t=\psi(x) \log x+O(x)
\end{align*}
$$

Inserting this estimate into (7.20), we have proved our first version of Selberg's fundamental formula: For $x \geq 1$,

$$
\begin{equation*}
\psi(x) \log x+\sum_{a b \leq x} \Lambda(a) \Lambda(b)=2 x \log x+O(x) \tag{7.22}
\end{equation*}
$$

It is convenient later to have a result expressed just in terms of primes and not prime powers. If we replace $\psi$ by $\theta$ on the left-hand side of (7.22), then we introduce an error of $\ll(\psi(x)-\theta(x)) \log x \ll x^{1 / 2}(\log x)^{2}$ (by (3.6)), which is certainly $O(x)$. Moreover, replacing

$$
\sum_{a b \leq x} \Lambda(a) \Lambda(b) \quad \text { by } \quad \sum_{p q \leq x} \log p \log q
$$

results in an error which is

$$
\begin{aligned}
& \ll \sum_{\substack{p^{a} q^{b} \leq x \\
a \geq 2 \text { or } b \geq 2}} \log p \log q \ll \sum_{\substack{p^{a} b^{b} \leq x \\
a \geq 2}} \log p \log q \ll \sum_{\substack{p^{a} \leq x \\
a \geq 2}} \log p \sum_{\substack{q^{b} \leq x / p^{a}}} \log q \\
& \ll \sum_{\substack{p^{a} \leq x \\
a \geq 2}}(\log p) \psi\left(x / p^{a}\right) \ll x \sum_{\substack{p^{a} \leq x \\
a \geq 2}} \frac{\log p}{p^{a}} \leq x \sum_{p} \frac{\log p}{p^{2}-p} \ll x,
\end{aligned}
$$

and this again fits within our existing error term. Thus

$$
\begin{equation*}
\theta(x) \log x+\sum_{p q \leq x} \log p \log q=2 x \log x+O(x) \tag{7.23}
\end{equation*}
$$

This is Selberg's formula in the shape (7.4) of the introduction, except that the second term in (7.4) appears as a sum over two variables here.

If we replace $\psi$ by $\theta$ in the calculation which gave (7.21), we find that $\sum_{p \leq x}(\log p)^{2}=\theta(x) \log x+O(x)$; this gives yet another form of the symmetry formula, which will also be helpful in the sequel:

$$
\begin{equation*}
\sum_{p \leq x}(\log p)^{2}+\sum_{p q \leq x} \log p \log q=2 x \log x+O(x) \tag{7.24}
\end{equation*}
$$

## 4. Removing the explicit appearance of primes

The goal of this section is to transition from the fundamental formula to the following inequality, where primes do not appear explicitly. Recall from
the introduction that the remainder-term function $R(x)$ is defined by the relation $\theta(x)=x+R(x)$.

Theorem 7.4. For $x \geq 1$, we have

$$
|R(x)| \log x \leq \int_{1}^{x}|R(x / t)| d t+O(x \log \log 3 x)
$$

The proof of Theorem 7.4 is not difficult, but it is somewhat long. We begin with a few routine but technical estimates.

Lemma 7.5. For $x \geq 1$, we have

$$
\sum_{p q \leq x} \frac{\log p \log q}{p q}=\frac{1}{2}(\log x)^{2}+O(\log x) .
$$

Proof. We have

$$
\begin{aligned}
\sum_{p q \leq x} \frac{\log p \log q}{p q} & =\sum_{p \leq x} \frac{\log p}{p} \sum_{q \leq x / p} \frac{\log q}{q}=\sum_{p \leq x} \frac{\log p}{p}(\log x-\log p+O(1)) \\
& =\log x \sum_{p \leq x} \frac{\log p}{p}-\sum_{p \leq x} \frac{(\log p)^{2}}{p}+O\left(\sum_{p \leq x} \frac{\log p}{p}\right) \\
& =\log x(\log x+O(1))-\sum_{p \leq x} \frac{(\log p)^{2}}{p}+O(\log x) \\
& =(\log x)^{2}-\sum_{p \leq x} \frac{(\log p)^{2}}{p}+O(\log x) .
\end{aligned}
$$

To handle the remaining sum we use partial summation. With $A(x)=$ $\sum_{p \leq x} p^{-1} \log p$, we have

$$
\begin{aligned}
\sum_{p \leq x} \frac{(\log p)^{2}}{p} & =A(x) \log x-\int_{1}^{x} \frac{A(t)}{t} d t \\
& =(\log x+O(1)) \log x-\int_{1}^{x} \frac{\log t+O(1)}{t} d t \\
& =(\log x)^{2}-\int_{1}^{x} \frac{\log t}{t} d t+O(\log x)=\frac{1}{2}(\log x)^{2}+O(\log x)
\end{aligned}
$$

inserting this estimate into (7.25) finishes the proof.
Lemma 7.6. For $x \geq 1$, we have

$$
\sum_{p q \leq x} \frac{\log p \log q}{p q \log (p q)}=\log x+O(\log \log 3 x)
$$

Proof. Let $a_{n}:=\sum_{p q=n} \frac{\log p \log q}{p q}$, and let $A(x):=\sum_{n \leq x} a_{n}$. Then $A(x)=$ $\frac{1}{2}(\log x)^{2}+O(\log x)$ for $x \geq 1$, by Lemma 7.5. So for $x \geq 3$, we have

$$
\begin{aligned}
\sum_{p q \leq x} \frac{\log p \log q}{p q \log (p q)} & =\frac{A(x)}{\log x}+\int_{2}^{x} \frac{A(t)}{t(\log t)^{2}} d t \\
& =\frac{1}{2} \log x+O(1)+\int_{2}^{x}\left(\frac{1}{2 t}+O\left(\frac{1}{t \log t}\right)\right) d t \\
& =\log x+O(\log \log x)
\end{aligned}
$$

Replacing $\log \log x$ by $\log \log 3 x$ ensures that the estimate is also valid for $1 \leq x \leq 3$.

Lemma 7.7. For $x \geq 1$, we have

$$
\sum_{p \leq x} \log p+\sum_{p q \leq x} \frac{\log p \log q}{\log (p q)}=2 x+O\left(\frac{x}{\log e x}\right) .
$$

Proof. With $A(x):=\sum_{p \leq x}(\log p)^{2}+\sum_{p q \leq x} \log p \log q$, the fundamental formula in the shape (7.24) supplies us with the estimate $A(x)=2 x \log x+$ $O(x)$. For $x \geq 2$, partial summation shows that

$$
\begin{aligned}
\sum_{p \leq x} \log p+\sum_{p q \leq x} & \frac{\log p \log q}{\log (p q)}=\frac{A(x)}{\log x}+\int_{2}^{x} \frac{A(t)}{t(\log t)^{2}} d t \\
& =2 x+O\left(\frac{x}{\log x}\right)+O\left(\int_{2}^{x} \frac{d t}{\log t}\right)=2 x+O\left(\frac{x}{\log x}\right)
\end{aligned}
$$

and this implies the stated result.
Lemma 7.8. For $x \geq 1$, we have

$$
\sum_{p q \leq x} \log p \log q=2 x \log x-\sum_{p q \leq x} \frac{\log p \log q}{\log (p q)} \theta(x / p q)+O(x \log \log 3 x)
$$

Proof. By Lemma 7.7 with $x$ replaced by $x / p$, we have

$$
\begin{aligned}
\sum_{p q \leq x} \log p \log q & =\sum_{p \leq x} \log p \sum_{q \leq x / p} \log q \\
& =\sum_{p \leq x} \log p\left(2 \frac{x}{p}-\sum_{q r \leq x / p} \frac{\log q \log r}{\log (q r)}+O\left(\frac{x / p}{\log e x / p}\right)\right) .
\end{aligned}
$$

This simplifies to

$$
\begin{aligned}
& 2 x \sum_{p \leq x} \frac{\log p}{p}-\sum_{p \leq x} \log p \sum_{q r \leq x / p} \frac{\log q \log r}{\log (q r)}+O\left(x \sum_{p \leq x} \frac{\log p}{p\left(1+\log \frac{x}{p}\right)}\right) \\
& =2 x \log x+O(x)-\sum_{q r \leq x} \frac{\log q \log r}{\log (q r)} \theta(x / q r)+O\left(x \sum_{p \leq x} \frac{\log p}{p\left(1+\log \frac{x}{p}\right)}\right) .
\end{aligned}
$$

To estimate the $O$-term, we partition those $p \leq x$ according to the integer $j \geq 0$ for which $e^{j} \leq x / p<e^{j+1}$; in this way we find

$$
\begin{aligned}
\sum_{p \leq x} \frac{\log p}{p\left(1+\log \frac{x}{p}\right)} & \leq \sum_{0 \leq j \leq \log x} \frac{1}{1+j} \sum_{x / e^{j+1} \leq p \leq x / e^{j}} \frac{\log p}{p} \\
& \ll \sum_{0 \leq j \leq \log x} \frac{1}{1+j} \ll \log \log 3 x
\end{aligned}
$$

Collecting these estimates and relabeling gives the statement of the lemma.

Lemma 7.9. For $x \geq 1$, we have

$$
\theta(x) \log x=\sum_{p q \leq x} \frac{\log p \log q}{\log (p q)} \theta\left(\frac{x}{p q}\right)+O(x \log \log 3 x) .
$$

Proof. According to Selberg's fundamental formula in the form (7.23), we have

$$
\theta(x) \log x=-\sum_{p q \leq x} \log p \log q+2 x \log x+O(x) .
$$

The result is obtained by replacing the right-hand sum with the estimate supplied for it by Lemma 7.8.

Proof of Theorem 7.4. We first re-express the fundamental formula as a relation involving $R(x)$. (Such a computation was alluded to in the introduction.) We have

$$
\begin{aligned}
R(x) \log x & =\theta(x) \log x-x \log x \\
& =\left(2 x \log x-\sum_{p q \leq x} \log p \log q\right)-x \log x+O(x) \\
& =x \log x-\sum_{p \leq x} \theta(x / p) \log p+O(x) .
\end{aligned}
$$

Replacing $\theta(x / p)$ with $x / p+R(x / p)$, we find that

$$
\begin{align*}
R(x) \log x & =x \log x-\sum_{p \leq x}\left(\frac{x}{p}+R\left(\frac{x}{p}\right)\right) \log p+O(x) \\
& =x \log x-x \sum_{p \leq x} \frac{\log p}{p}-\sum_{p \leq x} R\left(\frac{x}{p}\right) \log p+O(x) \\
& =-\sum_{p \leq x} R\left(\frac{x}{p}\right) \log p+O(x), \tag{7.26}
\end{align*}
$$

and so, in particular,

$$
\begin{equation*}
|R(x)| \log x \leq \sum_{p \leq x}|R(x / p)| \log p+O(x) . \tag{7.27}
\end{equation*}
$$

In order to deduce something like Theorem 7.4 from (7.27), we would like to have precise information about the partial sums of $\log p$. Of course such information is not available to us at this point! In order to work around this difficulty, we supplement (7.27) with another upper estimate on $|R(x)| \log x$ : By Lemma 7.9,

$$
\begin{aligned}
& R(x) \log x= \\
& \begin{aligned}
= & \theta(x) \log x-x \log x \\
p q \leq x & \frac{\log p \log q}{p q \log (p q)}+\sum_{p q \leq x} \frac{\log p \log q}{\log (p q)} R(x / p q) \\
& \quad x \log x+O(x \log \log 3 x) .
\end{aligned}
\end{aligned}
$$

Using Lemma 7.6 to estimate the first term here, we find that

$$
\begin{equation*}
R(x) \log x=\sum_{p q \leq x} \frac{\log p \log q}{\log (p q)} R(x / p q)+O(x \log \log 3 x) \tag{7.28}
\end{equation*}
$$

and so in particular,

$$
\begin{equation*}
|R(x)| \log x \leq \sum_{p q \leq x} \frac{\log p \log q}{\log (p q)}|R(x / p q)|+O(x \log \log 3 x) . \tag{7.29}
\end{equation*}
$$

Adding (7.27) to (7.29) shows that
(7.30) $\quad 2|R(x)| \log x$

$$
\leq \sum_{p \leq x} \log p|R(x / p)|+\sum_{p q \leq x} \frac{\log p \log q}{\log (p q)}|R(x / p q)|+O(x \log \log 3 x) .
$$

The contribution from the two sums on the right-hand side of (7.30) can be written in the form

$$
\sum_{n \leq x} a_{n}|R(x / n)|, \quad \text { where } \quad a_{n}:=\sum_{p=n} \log p+\sum_{p q=n} \frac{\log p \log q}{\log (p q)} .
$$

We are now in good shape, because we have an asymptotic formula for $A(x):=\sum_{n \leq x} a_{n}$; indeed, $A(x)=2 x+O(x / \log e x)$ by Lemma 7.7.

By Abel summation,

$$
\begin{aligned}
\sum_{n \leq x} a_{n}|R(x / n)| & =\sum_{n \leq x} A(n)\left|R\left(\frac{x}{n}\right)\right|-\sum_{n \leq x-1} A(n)\left|R\left(\frac{x}{n+1}\right)\right| \\
& =\sum_{n \leq x} A(n)\left(\left|R\left(\frac{x}{n}\right)\right|-\left|R\left(\frac{x}{n+1}\right)\right|\right)+O(x) .
\end{aligned}
$$

Substituting in our estimate for $A(x)$ and applying the triangle inequality, we deduce that the sum in (7.31) is

$$
\begin{align*}
2 \sum_{n \leq x} n\left(\left|R\left(\frac{x}{n}\right)\right|-\right. & \left.\left|R\left(\frac{x}{n+1}\right)\right|\right)  \tag{7.32}\\
& +O\left(\sum_{n \leq x} \frac{n}{1+\log n}\left|R\left(\frac{x}{n}\right)-R\left(\frac{x}{n+1}\right)\right|\right)
\end{align*}
$$

The main term of (7.32) telescopes to

$$
\begin{equation*}
2 \sum_{n \leq x}|R(x / n)|-2\lfloor x\rfloor R\left(\frac{x}{\lfloor x\rfloor+1}\right)=2 \sum_{n \leq x}|R(x / n)|+O(x) . \tag{7.33}
\end{equation*}
$$

To estimate the $O$-term in (7.32), we observe that

$$
\left|R\left(\frac{x}{n}\right)-R\left(\frac{x}{n+1}\right)\right|<\theta\left(\frac{x}{n}\right)-\theta\left(\frac{x}{n+1}\right)+\frac{x}{n^{2}},
$$

so that

$$
\begin{aligned}
\sum_{n \leq x} \frac{n}{1+\log n} & \left|R\left(\frac{x}{n}\right)-R\left(\frac{x}{n+1}\right)\right| \\
& \ll \sum_{n \leq x} \frac{n}{1+\log n}\left(\theta\left(\frac{x}{n}\right)-\theta\left(\frac{x}{n+1}\right)\right)+\sum_{n \leq x} \frac{x}{n(1+\log n)} .
\end{aligned}
$$

The latter sum on the right-hand side is $\ll x \log \log 3 x$, as we see by comparing with the corresponding integral. We rewrite the former sum, observing that

$$
\begin{aligned}
\sum_{n \leq x} \frac{n}{1+\log n} & \left(\theta\left(\frac{x}{n}\right)-\theta\left(\frac{x}{n+1}\right)\right) \\
& =\theta(x)+\sum_{n \leq x-1} \theta\left(\frac{x}{n+1}\right)\left(\frac{n+1}{1+\log (n+1)}-\frac{n}{1+\log n}\right) .
\end{aligned}
$$

Now $\theta(x) \ll x$. Moreover, since $\theta(x /(n+1)) \ll x / n$ and

$$
0 \leq \frac{n+1}{1+\log (n+1)}-\frac{n}{1+\log n} \leq \frac{1}{1+\log n}
$$

it follows that

$$
\begin{aligned}
\sum_{n \leq x-1} \theta\left(\frac{x}{n+1}\right)\left(\frac{n+1}{1+\log (n+1)}\right. & \left.-\frac{n}{1+\log n}\right) \\
& \ll x \sum_{n \leq x-1} \frac{1}{n(1+\log n)} \ll x \log \log 3 x .
\end{aligned}
$$

Collecting all of our estimates shows that

$$
\begin{equation*}
|R(x)| \log x \leq \sum_{n \leq x}|R(x / n)|+O(x \log \log 3 x) . \tag{7.34}
\end{equation*}
$$

In order to prove Theorem 7.4, we need to convert (7.34) into an inequality of integrals. To this end, observe that

$$
\begin{aligned}
\sum_{n \leq x}|R(x / n)|-\int_{1}^{x}|R(x / t)| d t & =\sum_{n \leq x} \int_{n}^{n+1}(|R(x / n)|-|R(x / t)|) d t+O(1) \\
& \leq \sum_{n \leq x} \int_{n}^{n+1}|R(x / n)-R(x / t)| d t+O(1)
\end{aligned}
$$

Now for $n \leq t \leq n+1$,

$$
\begin{aligned}
|R(x / n)-R(x / t)| & \leq \theta\left(\frac{x}{n}\right)-\theta\left(\frac{x}{t}\right)+\frac{x}{n}-\frac{x}{t} \\
& <\theta\left(\frac{x}{n}\right)-\theta\left(\frac{x}{n+1}\right)+\frac{x}{n^{2}}
\end{aligned}
$$

thus

$$
\begin{aligned}
\sum_{n \leq x}|R(x / n)|-\int_{1}^{x}|R(x / t)| d t & \leq \sum_{n \leq x}\left(\theta\left(\frac{x}{n}\right)-\theta\left(\frac{x}{n+1}\right)+\frac{x}{n^{2}}\right)+O(1) \\
& =\theta(x)+x \sum_{n \leq x} \frac{1}{n^{2}}+O(1) \ll x
\end{aligned}
$$

So by (7.34),

$$
|R(x)| \log x \leq \int_{1}^{x}|R(x / t)| d t+O(x \log \log 3 x)
$$

which is Theorem 7.4.

## 5. Nevanlinna's finishing strategy

5.1. Rescaling the remainder term. Put

$$
r(x):=e^{-x} R\left(e^{x}\right)=e^{-x} \theta\left(e^{x}\right)-1 .
$$

Our first goal is to prove the following analogue of Theorem 7.4 for $r(x)$, which appeared already in the introduction:

Theorem 7.10. As $x \rightarrow \infty$, we have

$$
\begin{equation*}
|r(x)| \leq \frac{1}{x} \int_{0}^{x}|r(t)| d t+o(1) \tag{7.35}
\end{equation*}
$$

Proof. We change variables in Theorem 7.4, replacing $t$ by $x / t$. This gives

$$
|R(x)| \log x \leq x \int_{1}^{x} \frac{|R(t)|}{t^{2}} d t+O(x \log \log 3 x)
$$

We now put $x=e^{y}$ and $t=e^{u}$ to find that

$$
\begin{aligned}
\left|R\left(e^{y}\right)\right| & \leq \frac{1}{y} e^{y} \int_{1}^{e^{y}} \frac{|R(t)|}{t^{2}} d t+O\left(e^{y} \frac{\log e y}{y}\right) \\
& =e^{y} \frac{1}{y} \int_{0}^{y}\left|R\left(e^{u}\right)\right| e^{-u} d u+O\left(e^{y} \frac{\log e y}{y}\right) .
\end{aligned}
$$

Theorem 7.10 follows upon multiplying both sides by $e^{-y}$.
The familiar estimate

$$
\sum_{p \leq x} \frac{\log p}{p}=\log x+O(1)
$$

also has a simple reformulation in terms of the function $r(x)$ :
Lemma 7.11. We have $\int_{0}^{x} r(t) d t=O(1)$ for all $x \geq 0$. As a consequence, there is a constant $C$ with

$$
\left|\int_{x}^{y} r(t) d t\right| \leq C
$$

for every pair of nonnegative real numbers $x$ and $y$.
Proof. We have

$$
\begin{aligned}
& \int_{0}^{x} r(t) d t=\int_{0}^{x} e^{-t} R\left(e^{t}\right) d t=\int_{1}^{e^{x}} \frac{R(u)}{u^{2}} d u=\int_{1}^{e^{x}} \frac{\theta(u)-u}{u^{2}} d u \\
& \quad=\int_{1}^{e^{x}} \frac{\theta(u)}{u^{2}} d u-\int_{1}^{e^{x}} \frac{d u}{u}=\left(\sum_{p \leq e^{x}} \frac{\log p}{p}-x\right)+O(1)=O(1) .
\end{aligned}
$$

5.2. Unbalancing the inequality. We say that $r$ changes sign at the point $x$ if there is a deleted neighborhood of $x$ on which $r$ has opposite signs to the left and right of $x$. Since $r$ is continuous except at the points of the set $\{\log p\}$ and is strictly decreasing between discontinuities, it is clear that $r$ has only countably many sign changes on $(0, \infty)$. Enumerate them as

$$
x_{1}<x_{2}<x_{3}<\cdots .
$$

(The ellipsis is not meant to imply that the sequence of $x_{i}$ is infinite; in fact it is - this follows from the work Littlewood alluded to in the notes to Chapter 3 - but this is immaterial for our purposes.) Whenever $x_{i}$ and $x_{i+1}$ are defined, let $I_{i}$ denote the half-open interval $\left[x_{i}, x_{i+1}\right)$.

Lemma 7.12. Suppose that $x$ and $x^{\prime}$ are consecutive sign changes of $r$, and let $I=\left[x, x^{\prime}\right)$.
(i) If $x$ is a sign change from negative to positive, then $r$ is positive on all of $I$. In this case $r$ is discontinuous at $x$.
(ii) If $x$ is a sign change from positive to negative, then $r$ is nonpositive on all of $I$. In this case $r$ is continuous at $x$.

Proof. Let $x$ be a change of sign from negative to positive. Since $r$ is strictly decreasing on each interval between discontinuities, it must be that $x=\log p_{0}$ for some prime $p_{0}$ and that $r(x)>0$. Suppose first that there are no primes $p$ with $x<\log p<x^{\prime}$. Then the restriction of $r$ to $I$ is continuous and strictly decreasing. This implies that any zero of $r$ on $I$ would be a change of sign in $r$. Since $x$ and $x^{\prime}$ are consecutive sign changes, $r$ is nonvanishing on $I$. So by continuity, $r$ is positive on all of $I$, as desired. If there are primes $p$ with $x<\log p<x^{\prime}$, then list the consecutive primes $p_{1}<\cdots<p_{k}$ for which

$$
\begin{equation*}
x<\log p_{1}<\log p_{2}<\cdots<\log p_{k}<x^{\prime} . \tag{7.36}
\end{equation*}
$$

The argument just given shows that there are no zeros of $r$ in $\left[x, \log p_{1}\right)$. So by continuity, $r$ is positive on $\left[x, \log p_{1}\right)$. Since $r$ has a positive jump at $\log p_{1}$, we have $r\left(\log p_{1}\right)>0$, and repeating the argument shows that $r$ is positive on $\left[\log p_{1}, \log p_{2}\right)$. Continuing in this way we eventually find that $r$ is positive on all of $I$.

Now suppose that $x$ is a sign change from positive to negative. Since $r$ has positive jumps at its discontinuities, $x$ must be a point of continuity of $r$, and so $r(x)=0$. If there are no primes $p$ with $x<\log p<x^{\prime}$, then $r$ is decreasing on $I$ and so the conclusion of (ii) is obvious. Otherwise, let $p_{1}<p_{2}<\cdots<p_{k}$ be the primes satisfying (7.36). We have that $r$ is negative on $\left(x, \log p_{1}\right)$, since $r(x)=0$ and $r$ is strictly decreasing on $\left[x, \log p_{1}\right)$. As a consequence, $r\left(\log p_{1}\right) \leq 0$, since otherwise $\log p_{1}$ would be a sign change between $x$ and $x^{\prime}$. Repeating this argument shows that
$r$ is nonpositive on each of the intervals $\left[\log p_{1}, \log p_{2}\right),\left[\log p_{2}, \log p_{3}\right), \ldots$, $\left[\log p_{k}, x^{\prime}\right)$; in fact, we find that $r$ is negative at each point of $I$ except possibly at $x$ and the $\log p_{i}$.

Put $\lambda:=\lim \sup |r(x)|$. The prime number theorem is the assertion that $\lambda=0$. Suppose for the sake of contradiction that $\lambda>0$, and fix a positive $\lambda^{\prime}$ with $\lambda^{\prime}<\min \{1, \lambda\}$. For each $i$ for which $I_{i}$ is defined, put

$$
I_{i}^{\prime}:=\left\{x \in I_{i}:|r(x)| \leq \lambda^{\prime}\right\} .
$$

The rest of this section is devoted to proving the following lemma:
Lemma 7.13. There is a constant $\kappa \in(0,1]$, depending only on $\lambda^{\prime}$, with $\mu\left(I_{i}^{\prime}\right) \geq \kappa \mu\left(I_{i}\right)$ whenever $I_{i}$ is defined.

Proof. Since $r$ does not change sign on $I_{i}=\left[x_{i}, x_{i+1}\right)$, Lemma 7.11 implies that

$$
\lambda^{\prime}\left(\mu\left(I_{i}\right)-\mu\left(I_{i}^{\prime}\right)\right)=\lambda^{\prime} \mu\left(I_{i} \backslash I_{i}^{\prime}\right) \leq \int_{I_{i} \backslash I_{i}^{\prime}}|r(t)| d t \leq\left|\int_{x_{i}}^{x_{i+1}} r(t) d t\right| \leq C,
$$

so that

$$
\begin{equation*}
\mu\left(I_{i}^{\prime}\right) \geq \mu\left(I_{i}\right)-C / \lambda^{\prime}=\left(1-\frac{C}{\lambda^{\prime} \mu\left(I_{i}\right)}\right) \mu\left(I_{i}\right) . \tag{7.37}
\end{equation*}
$$

This is enough to give the conclusion of Lemma 7.13 in the case when $\mu\left(I_{i}\right)$ is large (e.g., if $\mu\left(I_{i}\right) \geq 2 C / \lambda^{\prime}$ ). We now derive another estimate which will allow us to draw the same conclusion when $\mu\left(I_{i}\right)$ is small.

Lemma 7.12 implies that whenever $I_{i}$ is defined, precisely one of its endpoints is a point of continuity of $r$. Suppose it is the right endpoint $x_{i+1}$. Then $r\left(x_{i+1}\right)=0$, so that $\theta\left(e^{x_{i+1}}\right)=e^{x_{i+1}}$. Since $x_{i+1}$ represents a change from positive to negative, for each $x \in I_{i}$, we have

$$
0 \leq r(x)=e^{-x} \theta\left(e^{x}\right)-1 \leq e^{-x} \theta\left(e^{x_{i+1}}\right)-1=e^{x_{i+1}-x}-1 .
$$

In particular, $|r(x)| \leq \lambda^{\prime}$ for all $x \in I_{i}$ close enough to $x_{i+1}$, e.g., all $x \in I_{i}$ which satisfy

$$
x \geq x_{i+1}-\log \left(1+\lambda^{\prime}\right)
$$

(This situation is illustrated in Figure 1.) Similarly, if the left endpoint $x_{i}$ of $I_{i}$ is a point of continuity of $r$, then for each $x \in I_{i}$,

$$
0 \geq r(x)=e^{-x} \theta\left(e^{x}\right)-1 \geq e^{-x} \theta\left(e^{x_{i}}\right)-1=e^{x_{i}-x}-1
$$

Thus $|r(x)|=-r(x) \leq \lambda^{\prime}$ for all $x \in I_{i}$ close enough to $x_{i}$, say $x \leq x_{i}-$ $\log \left(1-\lambda^{\prime}\right)$. So certainly $|r(x)| \leq \lambda^{\prime}$ for those $x$ in the even smaller range $x \leq x_{i}+\log \left(1+\lambda^{\prime}\right)$.


Figure 1. Rough sketch of $r(x)$ vs. $e^{x_{i+1}-x}-1$ between the sign changes $x_{i}$ and $x_{i+1}$, in the case when $x_{i+1}$ is a point of continuity. Based on [Nev62].

So regardless of which endpoint of $I_{i}$ represents a point of continuity of $r$, we have

$$
\begin{align*}
\mu\left(I_{i}^{\prime}\right) & \geq \min \left\{\log \left(1+\lambda^{\prime}\right), \mu\left(I_{i}\right)\right\} \\
& =\min \left\{1, \log \left(1+\lambda^{\prime}\right) / \mu\left(I_{i}\right)\right\} \mu\left(I_{i}\right) . \tag{7.38}
\end{align*}
$$

Now (7.38) is the sought-after dual to (7.37); together these estimates imply Lemma 7.13. Indeed, set $C^{\prime}:=2 C / \lambda^{\prime}$. If $\mu\left(I_{i}\right) \geq C^{\prime}$, then we have $\mu\left(I_{i}^{\prime}\right) \geq$ $\frac{1}{2} \mu\left(I_{i}\right)$ by (7.37). Otherwise

$$
\mu\left(I_{i}^{\prime}\right) \geq \min \left\{1, \log \left(1+\lambda^{\prime}\right) / C^{\prime}\right\} \mu\left(I_{i}\right)
$$

by (7.38). So Lemma 7.13 follows with $\kappa=\min \left\{1 / 2, \log \left(1+\lambda^{\prime}\right) / C^{\prime}\right\}$.
5.3. Endgame. We can now complete the proof of the prime number theorem in the form $r(x)=o(1)$. Let $\epsilon>0$ be arbitrary but fixed. By the choice of $\lambda$, we can select a positive number $x_{\epsilon}$ with the property that for every $x>x_{\epsilon}$, we have

$$
|r(x)| \leq \lambda+\epsilon
$$

Now let $x$ be large, and let $x_{\epsilon}<x_{m}<x_{m+1}<\cdots<x_{n} \leq x$ be a list of the sign changes of $r$ in $\left(x_{\epsilon}, x\right]$. If there are no sign changes in this interval, then

$$
\int_{0}^{x}|r(t)| d t \leq \int_{0}^{x_{\epsilon}}|r(t)| d t+\left|\int_{x_{\epsilon}}^{x} r(t) d t\right|=O(1)
$$

by Lemma 7.11. Otherwise, by Lemmas 7.11 and 7.13, we have

$$
\begin{align*}
\int_{0}^{x}|r(t)| d t & \leq \int_{0}^{x_{m}}|r(t)| d t+\sum_{j=m}^{n-1} \int_{x_{j}}^{x_{j+1}}|r(t)| d t+\int_{x_{n}}^{x}|r(t)| d t \\
& \leq \sum_{j=m}^{n-1}\left((\lambda+\epsilon)\left(\mu\left(I_{j}\right)-\mu\left(I_{j}^{\prime}\right)\right)+\lambda^{\prime} \mu\left(I_{j}^{\prime}\right)\right)+O(1) \\
& =\sum_{j=m}^{n-1}\left((\lambda+\epsilon) \mu\left(I_{j}\right)+\left(\lambda^{\prime}-\lambda-\epsilon\right) \mu\left(I_{j}^{\prime}\right)\right)+O(1) \\
& \leq \sum_{j=m}^{n-1}\left((\lambda+\epsilon) \mu\left(I_{j}\right)+\kappa\left(\lambda^{\prime}-\lambda-\epsilon\right) \mu\left(I_{j}\right)\right)+O(1) \\
& \leq\left((\lambda+\epsilon)+\kappa\left(\lambda^{\prime}-\lambda-\epsilon\right)\right) x+O(1) . \tag{7.39}
\end{align*}
$$

Therefore, by Theorem 7.10,

$$
\limsup _{x \rightarrow \infty}|r(x)| \leq \limsup _{x \rightarrow \infty} \frac{1}{x} \int_{0}^{x}|r(t)| d t \leq(\lambda+\epsilon)+\kappa\left(\lambda^{\prime}-\lambda-\epsilon\right) .
$$

Since this holds for each $\epsilon>0$, letting $\epsilon \downarrow 0$, it follows that

$$
\limsup _{x \rightarrow \infty}|r(x)| \leq \lambda+\kappa\left(\lambda^{\prime}-\lambda\right)<\lambda,
$$

contradicting that $\lambda=\lim \sup |r(x)|$. We have proved the prime number theorem!

## Notes

In addition to the original papers of Nevanlinna and Selberg, our organization of the proof of the prime number theorem has been influenced heavily by Shapiro's treatment [Sha83, Chapter 10]. Our account of the early history of the elementary proof of the prime number theorem is based on the recollections of Selberg (as recorded in [Gol04], [BS08]) and Straus [Str].

The Wiener-Ikehara theorem and its application to the prime number theorem are described, for example, in [Mur01, Chapter 3]. The approach to the PNT via theorems of this type (so-called "Tauberian theorems") is discussed extensively in [Kor02, $\S \S 1-8]$; see also $[\mathbf{N a r 0 4}, \S 6.4]$. Using little more than Cauchy's integral theorem, one can prove a weak version of the Wiener-Ikehara result that suffices for the proof of the prime number theorem. In this way one obtains the simplest known analytic proof. The
groundwork for these developments was laid by Newman ([New80], see also [New98, Chapter VII]). Polished versions of this argument appear in papers of Korevaar [Kor82] and Zagier [Zag97], and a very readable account of the method is given in the text of Hlawka et al. [HST91].

The Erdős-Selberg method applies also to certain generalizations of the prime number theorem. In particular, their method leads to an elementary proof of the prime number theorem for arithmetic progressions (see [Sel50]) as well as a proof of the "prime ideal theorem" of algebraic number theory (see [Sha49a]). One version of the argument for arithmetic progressions is sketched in the exercises (cf. [Nev64, Kapitel III]).

Until 1980 all extant elementary proofs of the prime number theorem were variants on the Erdős-Selberg argument, at their core relying on some version of Selberg's fundamental formula. Since then Daboussi [Dab84] and Hildebrand [Hil86] have given proofs independent of the Erdős-Selberg work. Daboussi's argument is described at length in the engaging monograph of Tenenbaum \& Mendès France [TMF00].

So far the elementary proofs of the prime number theorem have not had the dramatic effect on the number-theoretic landscape predicted by Hardy. Rather than overthrow the use of complex-variable methods, the existing elementary proofs have shown themselves to be comparatively inflexible, and if anything have underscored the utility of analytic techniques. For example, no elementary proof of the prime number theorem is known which gives an estimate for the error term of the same quality as what was obtained by de la Vallée-Poussin already in 1899 (cf. the notes to Chapter 3).

## Exercises

Analytic exercises. This series of exercises requires familiarity with complex analysis. Unless otherwise specified, $s$ denotes a complex variable and we write $s=\sigma+i \tau$, where $\sigma, \tau \in \mathbf{R}$.

1. Prove (7.2) by computing $\int t^{-s} d A(t)$ for $A(x):=\sum_{n \leq x} 1$.
2. (Dirichlet-Dedekind [Dir99, §118]) Suppose that $f$ is a complex-valued arithmetic function whose partial sums satisfy

$$
\sum_{n \leq x} f(n)=(R+o(1)) x
$$

for some complex number $R$ (as $x \rightarrow \infty$ ). Prove that the Dirichlet series $F(s):=\sum_{n=1}^{\infty} f(n) / n^{s}$ converges to a continuous function on $(1, \infty)$ and that

$$
\lim _{s \downarrow 1}(s-1) F(s)=R .
$$

(Theorem 7.1 may be viewed as a sort of converse of this result.)
3. (Chebyshev [Che51]; cf. [Nar04, pp. 100-102]) In Chapter 3 we proved the theorem of Chebyshev (Theorem 3.4) that if $\frac{\pi(x)}{x / \log x}$ tends to a limit as $x \rightarrow \infty$, then that limit is necessarily 1 . Actually Chebyshev proved a stronger result: If we put $E(x):=\pi(x)-\operatorname{Li}(x)$, then for each natural number $k$,

$$
\begin{equation*}
\limsup _{x \rightarrow \infty} \frac{E(x)}{x /(\log x)^{k}} \geq 0 \quad \text { and } \quad \liminf _{x \rightarrow \infty} \frac{E(x)}{x /(\log x)^{k}} \leq 0 . \tag{7.40}
\end{equation*}
$$

In this exercise we sketch a proof of this result.
Let $k$ be a natural number. For real $s>1$, put $P(s):=\sum_{p} p^{-s}$.
(a) Show that $P(s)-\log \frac{1}{s-1}$ has an analytic continuation to an open subset of the complex plane containing all real $s \geq 1$. Deduce that if we put

$$
F(s):=(-1)^{k}\left(P^{(k)}(s)-\zeta^{(k-1)}(s)\right),
$$

then $F(s)$ remains bounded as $s$ tends to 1 from above.
(b) Show that for $s>1$, we have $F(s)=\sum_{p} \frac{(\log p)^{k}}{p^{s}}-\sum_{n=2}^{\infty} \frac{(\log n)^{k-1}}{n^{s}}$.
(c) Show that for $s>1$,

$$
F(s)=-\int_{2}^{\infty} E(t) \frac{d}{d t}\left(\frac{(\log t)^{k}}{t^{s}}\right) d t+O_{k}(1) .
$$

(d) Deduce (7.40) from (a) and (c). Check that when $k=1$, these inequalities imply Theorem 3.4.

Remark. It follows from Exercise 3.8 (which assumes the prime number theorem with a reasonable error term) that both limits in (7.40) vanish for each $k$.
4. Define $Z(s):=-\frac{\zeta^{\prime}(s)}{\zeta(s)}$. From (7.3) we know that

$$
Z(s)=\sum_{n=1}^{\infty} \frac{\Lambda(n)}{n^{s}} \quad(\sigma>1)
$$

(a) Prove that for $\sigma>1$,

$$
Z(s)=\frac{s}{s-1}+s \int_{1}^{\infty}(\psi(t)-t) \frac{d t}{t^{s+1}} .
$$

(b) Assuming the prime number theorem in the form $\psi(x) \sim x$, show that the integral in (a) is $o(1 /(\sigma-1))$, as $\sigma \downarrow 1$, uniformly in $\tau$. Conclude that for fixed $\tau \neq 0$, one has

$$
\lim _{\sigma \downarrow 1}(\sigma-1)|Z(\sigma+i \tau)|=0 .
$$

(c) On the other hand, show that if $\zeta(s)$ has a zero of order $m \geq 0$ at $1+i \tau$ (so that necessarily $\tau \neq 0$ ), then

$$
\lim _{\sigma \downarrow 1}(\sigma-1) Z(\sigma+i \tau)=-m
$$

Combining the results of (b) and (c), deduce that if the prime number theorem is true, then $\zeta(s)$ has no zeros on the line $\sigma=1$.
5. Let $M(x):=\sum_{n \leq x} \mu(n)$, where $\mu$ is the Möbius function.
(a) Prove that if $M(x) / x$ tends to a limit, then that limit must be zero. Hint: Use Corollary 7.3(i) or the result of Exercise 2.
(b) Assuming that $\zeta(s)$ has no zeros on the line $\sigma=1$, deduce from the Wiener-Ikehara Theorem (Theorem 7.1) that $M(x) / x$ does, in fact, tend to zero.
(c) Suppose, conversely, that $M(x) / x$ tends to zero. Prove that $\zeta(s)$ has no zeros on the line $\Re(s)=1$. (Cf. Exercise 4.)
The estimate $M(x)=o(x)$ can be interpreted probabilistically: If a squarefree number $n$ is chosen at random, it is equally likely to have an even number of prime factors as an odd number of prime factors.

Remark. From (b) and (c), we see that the estimate $M(x)=o(x)$ is in some sense equivalent to the prime number theorem, since both amount to the nonexistence of zeros of $\zeta(s)$ on the line $\sigma=1$. For an elementary proof of this equivalence, see [Apo76, §4.9].
6. We outline a proof, taken from [Tit86, §3.2], that $\zeta(s)$ is nonvanishing on the line $\sigma=1$. We assume the result of Exercise 1, so that $\zeta(s)$ is
known to be analytic for $\sigma>0$ except for a simple pole at $s=1$ with residue 1. As before we let $Z(s)=-\frac{\zeta^{\prime}(s)}{\zeta(s)}$.
(a) Suppose that $\zeta(s)$ has a zero at $s_{0}=1+i \tau_{0}$, where necessarily $\tau_{0} \neq 0$. Prove that $s_{0}$ is necessarily simple. Hint: If $\zeta(s)$ has a zero of order $k$ at $s_{0}$, then $Z(s)$ has a simple pole at $s_{0}$ with residue $-k$; however,

$$
\left|\frac{\zeta^{\prime}\left(\sigma+i \tau_{0}\right)}{\zeta\left(\sigma+i \tau_{0}\right)}\right| \leq\left|\sum_{n=1}^{\infty} \frac{\Lambda(n)}{n^{\sigma+i \tau_{0}}}\right| \leq \sum_{n=1}^{\infty} \frac{\Lambda(n)}{n^{\sigma}} \sim \frac{1}{\sigma-1} \quad \text { as } \sigma \downarrow 1 .
$$

(b) Show that

$$
\sum_{n=1}^{\infty} \frac{\Lambda(n) \cos \left(\tau_{0} \log n\right)}{n^{\sigma}} \sim-\frac{1}{\sigma-1} \quad \text { as } \sigma \downarrow 1
$$

(c) By the Cauchy-Schwarz inequality, we now have

$$
\begin{aligned}
\frac{1}{(\sigma-1)^{2}} & \sim\left(\sum_{n=1}^{\infty} \frac{\Lambda(n) \cos \left(\tau_{0} \log n\right)}{n^{\sigma}}\right)^{2} \\
& \leq\left(\sum_{n=1}^{\infty} \frac{\Lambda(n) \cos ^{2}\left(\tau_{0} \log n\right)}{n^{\sigma}}\right)\left(\sum_{n=1}^{\infty} \frac{\Lambda(n)}{n^{\sigma}}\right)
\end{aligned}
$$

Rewriting $\cos ^{2}\left(\tau_{0} \log n\right)=\frac{1}{2}\left(1+\cos \left(2 \tau_{0} \log n\right)\right)$ and using that $Z(s)$ has a simple pole at $s=1$ with residue 1 , prove that

$$
\sum_{n=1}^{\infty} \frac{\Lambda(n) \cos \left(2 \tau_{0} \log n\right)}{n^{\sigma}} \geq(1+o(1)) \frac{1}{\sigma-1} \quad \text { as } \sigma \downarrow 1 \text {. }
$$

(d) Show that the estimate of (c) contradicts the regularity of $\zeta(s)$ at the point $1+2 i \tau_{0}$.

In Exercises 7-10, we look a bit deeper at how an elementary proof of the prime number theorem is possible given its equivalence to the nonvanishing of $\zeta(s)$ on the line $\sigma=1$. This paradox was addressed by Ingham in his expert review [Ing48] of the Erdős-Selberg papers. Following Ingham, suppose that $F$ is any function of a complex variable with the following three properties:
(i) $F$ is represented by the Dirichlet series $\sum_{n=1}^{\infty} a_{n} / n^{s}$ in the region $\sigma>1$, where each $a_{n}$ is real.
(ii) $F$ is analytic on the closed half-plane $\sigma \geq 1$, except possibly for simple poles on the line $\sigma=1$.
(iii) For $\sigma>1$, we have $-F^{\prime}(s)+F(s)^{2}=\sum_{n=1}^{\infty} b_{n} / n^{s}$, where $B(x):=$ $\sum_{n \leq x} b_{n}$ satisfies $B(x) \sim 2 x \log x$ as $x \rightarrow \infty$.

The reader should have in the back of her mind the special case $F(s)=$ $-\frac{\zeta^{\prime}(s)}{\zeta(s)}$. Then (i) holds with $a_{n}=\Lambda(n)$, and (iii) is a consequence of Selberg's fundamental formula in the shape (7.20).
7. Let $G(s):=-F^{\prime}(s)+F(s)^{2}+2 \zeta^{\prime}(s)$. Show that $G$ is represented by a Dirichlet series $\sum c_{n} / n^{s}$ where $\sum_{n \leq x} c_{n}=o(x \log x)$. Deduce that $G$ has no poles of order $\geq 2$ on $\sigma=1$, and hence that $-F^{\prime}(s)+F(s)^{2}$ has no poles of order $\geq 2$ on $\sigma=1$ except possibly at $s=1$.
8. Deduce from Exercise 7 that if $F$ has a pole at $1+i \tau_{0}$ with $\tau_{0} \neq 0$, then its residue $R$ there satisfies $R+R^{2}=0$. Conclude that $R=-1$.
9. We now describe how to construct a function $F$ possessing properties (i)-(iii) which nevertheless has a pole on the line $\sigma=1$ other than $s=1$.
(a) For each fixed real number $\alpha \neq 0$, show that

$$
\sum_{n \leq x} n^{i \alpha}=\frac{1}{1+i \alpha} x^{1+i \alpha}+o(x) \quad \text { while } \quad \sum_{n \leq x} \frac{1}{n} n^{i \alpha}=o(\log x) .
$$

(b) Show that $F(s):=\zeta(s)-\zeta(s-i \alpha)-\zeta(s+i \alpha)$ possesses each of Ingham's properties (i)-(iii). Of course (i) and (ii) are immediate; establishing (iii) is the difficult component and where the estimates of (a) come into play.
(c) Show that $F$ has a pole at $s=1+i \alpha$.

Exercise 9 shows that (i)-(iii) are not enough to rule out poles of $F$ of the form $1+i \tau_{0}$, with $\tau_{0} \neq 0$. Quoting Ingham (ibid.),
this may be taken as a reason why it is possible to give an elementary proof of [Selberg's fundamental formula] without becoming involved in the question of the existence of zeros of $\zeta$ on $\sigma=1$.
10. Now suppose that in addition to (i)-(iii) we require that each $a_{n} \geq 0$, i.e., that $F$ is represented by a Dirichlet series with nonnegative coefficients. (This is not satisfied for the $F$ of Exercise 9.) If $F$ has a pole at $1+i \tau_{0}$ with $\tau_{0} \neq 0$, then by assumption (ii) and Exercise 8, this pole is simple with residue -1 . By imitating the argument of Exercise 6, show that this forces $F$ to have a pole of residue $\geq 1$ at $1+2 i \tau_{0}$, contradicting Exercise 8.
Taking $F(s)=-\frac{\zeta^{\prime}(s)}{\zeta(s)}$ in Exercise 10, we see that the nonvanishing of $\zeta(s)$ on $\sigma=1$ is a consequence of Selberg's fundamental formula paired with the nonnegativity of $\Lambda(n)$.

Primes in progressions. Recall that $\pi(x ; m, a)$ denotes the number of primes $p \leq x$ with $p \equiv a(\bmod m)$. The next series of exercises leads the
reader through a proof of the following fundamental equidistribution result, already alluded to in Chapter 1.
$\star$ Theorem 7.14 (The prime number theorem for arithmetic progressions). Suppose that $a$ and $m$ are relatively prime integers with $m>0$. Then

$$
\pi(x ; m, a) \sim \frac{1}{\varphi(m)} \frac{x}{\log x} \quad(x \rightarrow \infty) .
$$

The steps in the proof of Theorem 7.14 correspond closely to those in the proof of the prime number theorem (which is the case $m=1$ of Theorem 7.14). However, in place of Mertens' estimate for the partial sums of $\log p / p$, we make frequent use of the deeper result that for $(a, m)=1$,

$$
\begin{equation*}
\sum_{\substack{p \leq x \\ p \equiv a \leq \bmod m}} \frac{\log p}{p}=\frac{1}{\varphi(m)} \log x+O(1), \tag{7.41}
\end{equation*}
$$

which we established elementarily in the course of proving Dirichlet's theorem. In (7.41) and the exercises below, all the implied constants are allowed to depend on $m$.

Define

$$
\theta(x ; m, a):=\sum_{\substack{p \leq x \\ p \equiv a(\bmod m)}} \log p .
$$

11. Let $a$ and $m$ be coprime integers with $m>0$. Prove that as $x \rightarrow \infty$,

$$
\pi(x ; m, a) \sim \frac{1}{\varphi(m)} \frac{x}{\log x} \Longleftrightarrow \theta(x ; m, a) \sim \frac{x}{\varphi(m)} .
$$

It is in this latter form that Theorem 7.14 will be established.
We begin the demonstration of Theorem 7.14 by proving an analogue of Selberg's fundamental formula:

* Theorem 7.15 (Selberg's formula for arithmetic progressions). Let $a$ and $m$ be coprime integers with $m>0$. Then for $x \geq 1$,

$$
\begin{equation*}
\sum_{\substack{p \leq x \\ p \equiv a(\bmod m)}}(\log p)^{2}+\sum_{\substack{p q \leq x \\ p q \equiv a(\bmod m)}} \log p \log q=\frac{2}{\varphi(m)} x \log x+O(x) . \tag{7.42}
\end{equation*}
$$

12. (a) Fix a coprime residue class $a \bmod m$. By summing both sides of the identity (7.17) over the progression $a \bmod m$, show that

$$
\begin{gather*}
\sum_{\substack{n \leq x \\
n \equiv a(\bmod m)}} \Lambda(n) \log n+\sum_{\substack{d f \leq x \\
d f \equiv a(\bmod m)}} \Lambda(d) \Lambda(f)  \tag{7.43}\\
=\sum_{\substack{d \leq x \\
(d, m)=1}} \mu(d) \sum_{\substack{f \leq x / d \\
f \equiv a \bar{d}(\bmod m)}}(\log f)^{2} \\
=\frac{1}{m} \sum_{\substack{d f \leq x \\
(d f, m)=1}} \mu(d)(\log f)^{2}+O(x),
\end{gather*}
$$

where $\bar{d}$ denotes a solution of $d \bar{d} \equiv 1(\bmod m)$.
(b) Summing over all invertible residue classes $a \bmod m$, deduce that

$$
\sum_{n \leq x} \Lambda(n) \log n+\sum_{d f \leq x} \Lambda(d) \Lambda(f)=\frac{\varphi(m)}{m} \sum_{\substack{d f \leq x \\(d f, m)=1}} \mu(d)(\log f)^{2}+O(x)
$$

The left-hand side here coincides with that of Selberg's original fundamental formula (in the form (7.20)). Deduce from that formula and (7.43) that

$$
\begin{equation*}
\sum_{\substack{n \leq x \\ n \equiv a(\bmod m)}} \Lambda(n) \log n+\sum_{\substack{d f \leq x \\ d f \equiv a(\bmod m)}} \Lambda(d) \Lambda(f)=\frac{2}{\varphi(m)} x \log x+O(x) . \tag{7.44}
\end{equation*}
$$

(c) Deduce Theorem 7.15 from (7.44) by showing that the contribution in (7.44) from proper prime powers is $O(x)$.
Define the remainder term $R(x ; m, a)$ by

$$
R(x ; m, a):=\theta(x ; m, a)-\frac{x}{\varphi(m)} .
$$

Theorem 7.14 amounts to the assertion that $R(x ; m, a)=o(x)$ whenever $a$ and $m$ are coprime integers with $m>0$.

To proceed we need the following analogue of Theorem 7.4: If $a$ and $m$ are coprime integers with $m>0$, then for $x \geq 1$,

$$
\begin{equation*}
|R(x ; m, a)| \log x \leq \frac{1}{\varphi(m)} \sum_{\substack{b \text { mod } m \\(b, m)=1}} \int_{1}^{x}|R(x / t ; m, b)| d t+O(x \log \log 3 x) . \tag{7.45}
\end{equation*}
$$

13. Here is an outline of the proof of (7.45). Let $a$ and $m$ be coprime integers with $m>0$. Prove that each of the following estimates holds for all $x \geq 1$ :
(a) (Cf. Lemma 7.5)

$$
\sum_{\substack{p q \leq x \\=a(\bmod m)}} \frac{\log p \log q}{p q}=\frac{1}{2 \varphi(m)}(\log x)^{2}+O(\log x)
$$

(b) (Cf. Lemma 7.6)

$$
\sum_{\substack{p q \leq x \\=a(\bmod m)}} \frac{\log p \log q}{p q \log (p q)}=\frac{1}{\varphi(m)} \log x+O(\log \log 3 x)
$$

(c) (Cf. Lemma 7.7)

$$
\sum_{\substack{p \leq x \\ a(\bmod m)}} \log p+\sum_{\substack{p q \leq x \\ p q \equiv a(\bmod m)}} \frac{\log p \log q}{\log (p q)}=\frac{2}{\varphi(m)} x+O\left(\frac{x}{\log e x}\right)
$$

(d) (Cf. Lemma 7.8)

$$
\begin{array}{r}
\sum_{\substack{p q \leq x \\
p q \equiv a(\bmod m)}} \log p \log q=\frac{2}{\varphi(m)} x \log x-\sum_{\substack{p q \leq x \\
(p q, m)=1}} \frac{\log p \log q}{\log (p q)} \theta(x / p q ; m, a \overline{p q}) \\
+O(x \log \log 3 x),
\end{array}
$$

where $\overline{p q}$ is an inverse of $p q$ modulo $m$.
(e) (Cf. Lemma 7.9)

$$
\begin{aligned}
\theta(x ; m, a) \log x=\sum_{\substack{p q \leq x \\
p q \equiv a(\bmod m)}} \frac{\log p \log q}{\log (p q)} \theta(x / p q ; m, a \overline{p q}) & \\
& +O(x \log \log 3 x) .
\end{aligned}
$$

(f) (Cf. (7.26))

$$
R(x ; m, a) \log x=-\sum_{\substack{p \leq x \\ p \nmid m}} R(x / p ; m, a \bar{p}) \log p+O(x) .
$$

(g) (Cf. (7.28))

$$
R(x ; m, a) \log x=\sum_{\substack{p q \leq x \\(p q, m)=1}} \frac{\log p \log q}{\log (p q)} R(x / p q ; m, a \overline{p q})+O(x \log \log 3 x) .
$$

(h) Suppose that $b$ is coprime to $m$. With $\bar{b}$ denoting an inverse of $b$ modulo $p$, show that

$$
\begin{array}{r}
\sum_{\substack{p \leq x \\
p \equiv a \bar{b}(\bmod m)}}|R(x / p ; m, b)| \log p+\sum_{\substack{p q \leq x \\
p q \equiv a \bar{b}(\bmod m)}} \frac{\log p \log q}{\log (p q)}|R(x / p q ; m, b)| \\
=\frac{2}{\varphi(m)} \sum_{n \leq x}|R(x / n ; m, b)|+O(x \log \log 3 x) .
\end{array}
$$

(i) Combining the results of (f)-(h), prove that

$$
\begin{equation*}
|R(x ; m, a)| \log x \leq \frac{1}{\varphi(m)} \sum_{\substack{b \text { mod } \\(b, m)=1}} \sum_{n \leq x}|R(x / n ; m, b)|+O(x \log \log 3 x) . \tag{7.46}
\end{equation*}
$$

(j) By replacing the inner sum in (7.46) by the corresponding integral, prove the original claim (7.45).
14. We can now complete the proof of the prime number theorem for arithmetic progressions (Theorem 7.14). Define a rescaled remainder term function $r(x ; m, a)$ by

$$
\begin{aligned}
r(x ; m, a) & :=e^{-x} R\left(e^{x} ; m, a\right) \\
& =e^{-x} \theta\left(e^{x} ; m, a\right)-\frac{1}{\varphi(m)} .
\end{aligned}
$$

Fix a positive integer $m$. The prime number theorem for primes in residue classes modulo $m$ is the assertion that $r(x ; m, a)=o(1)$ whenever $\operatorname{gcd}(a, m)=1$. Put

$$
\lambda:=\max _{(a, m)=1} \limsup |r(x ; m, a)| .
$$

Assume for the sake of contradiction that $\lambda>0$.
(a) (Cf. Theorem 7.10) Show that if $a$ is coprime to $m$, then as $x \rightarrow \infty$,

$$
|r(x ; m, a)| \leq \frac{1}{\varphi(m)} \sum_{\substack{b \bmod m \\(b, m)=1}} \frac{1}{x} \int_{0}^{x}|r(t ; m, b)| d t+o(1) .
$$

(b) (Cf. Lemma 7.11) Prove that if $a$ is coprime to $m$, then

$$
\left|\int_{x}^{y} r(t ; m, a) d t\right| \leq C
$$

for all nonnegative real numbers $x$ and $y$, where $C$ is a constant depending only on $m$.
(c) By mimicking the arguments of $\S \S 5.2-5.3$, show that whenever $b$ is coprime to $m$,

$$
\limsup _{x \rightarrow \infty} \frac{1}{x} \int_{0}^{x}|r(t ; m, b)| d t<\lambda .
$$

Deduce from part (a) that whenever $a$ is coprime to $m$,

$$
\lim \sup |r(x ; m, a)|<\lambda
$$

Since this holds for each $a$ coprime to $m$, this contradicts the definition of $\lambda$. Hence $\lambda=0$.
This completes the proof of Theorem 7.14.

## Miscellany.

15. Let $\beta(n)=\sum_{p \mid n} p$ be the sum of the distinct prime divisors of $n$ and let $p(n)$ and $P(n)$ denote the smallest and largest prime factors of $n$. Show that as $x \rightarrow \infty$,

$$
\begin{aligned}
\sum_{n \leq x} \beta(n) \sim \sum_{2 \leq n \leq x} P(n) & \sim \frac{\zeta(2)}{2} \frac{x^{2}}{\log x}, \text { and } \\
\sum_{p \leq x} p & \sim \sum_{2 \leq n \leq x} p(n) \sim \frac{1}{2} \frac{x^{2}}{\log x} .
\end{aligned}
$$

The results for $\beta(n)$ and $p(n)$ here are due to Kalecki [Kal64]; a sharper form of the $P(n)$ estimate is due to Brouwer [Bro74].
16. (Moser [Mos63]; see also [Guy04, C2]) Let $r(n)$ be the number of ways of writing the natural number $n$ as a sum of consecutive primes. For example, $r(83)=3$, since 83 has the three representations

$$
11+13+17+19+23, \quad 23+29+31, \quad 83
$$

Show that $r(n)$ has mean value $\log 2$; in other words, prove that as $x \rightarrow \infty$,

$$
\frac{1}{x} \sum_{n \leq x} r(n) \rightarrow \log 2 .
$$

Hint: For each natural number $k$, let $r(n, k)$ be the number of ways of writing $n$ as a sum of $k$ consecutive primes (so $r(n, k)$ is 0 or 1 for each $n)$. Begin by showing that $\pi(x / k)-(k-1) \leq \sum_{n \leq x} r(n, k) \leq \pi(x / k)$.
17. (Cf. Mirsky [Mir49]) Let $\pi^{*}(x)$ denote the number of primes $p \leq x$ for which $p-1$ is squarefree. In this exercise we outline a proof, based on the prime number theorem for arithmetic progressions (Theorem 7.14) and the Brun-Titchmarsh inequality (Exercise 6.21), that as $x \rightarrow \infty$,

$$
\pi^{*}(x) \sim A \frac{x}{\log x} \quad \text { where } \quad A=\prod_{q}\left(1-\frac{1}{q(q-1)}\right) .
$$

(Here, as usual, $q$ denotes a prime variable.) In particular, $p-1$ is squarefree for a positive proportion of primes $p$. The constant of proportionality $A$ is known as Artin's constant and $A \approx 0.3739558$.
(a) Let $z>0$ be arbitrary but fixed. Using Theorem 7.14 and the principle of inclusion-exclusion, show that as $x \rightarrow \infty$,

$$
\sum_{\substack{p \leq x \\ q^{2} \mid p-1 \Rightarrow q>z}} 1 \sim A_{z} \frac{x}{\log x}, \quad \text { where } \quad A_{z}=\prod_{q \leq z}\left(1-\frac{1}{q(q-1)}\right) .
$$

(b) The sum in (a) majorizes $\pi^{*}(x)$, since it includes every prime $p$ for which $p-1$ is squarefree. On the other hand, if $p$ is counted by that sum but $p-1$ is not squarefree, then $p-1$ is divisible by $q^{2}$ for some prime $q>z$. It follows that

$$
\sum_{\substack{p \leq x \\ \mid p-1 \Rightarrow q>z}} 1 \leq \pi^{*}(x)+\sum_{q>z} \pi\left(x ; q^{2}, 1\right)
$$

Using the Brun-Titchmarsh inequality, show that the terms of the right-hand sum corresponding to $q \in\left(z,(\log x)^{2}\right]$ contribute $\ll x /(z \log x)$. Using the trivial bound $\pi\left(x ; q^{2}, 1\right) \leq x / q^{2}$, show that the terms with $q>(\log x)^{2}$ contribute $\ll x /(\log x)^{2}$.
(c) Deduce from (a) and (b) that for each fixed $z$, $\limsup _{x \rightarrow \infty} \frac{\pi^{*}(x)}{x / \log x} \leq A_{z} \quad$ while $\quad \liminf _{x \rightarrow \infty} \frac{\pi^{*}(x)}{x / \log x} \geq A_{z}-O(1 / z)$. Complete the proof by letting $z \rightarrow \infty$.

## Perfect Numbers and their Friends


#### Abstract

Among all the problems which we are used to dealing with in Mathematics, none for certain, are judged by the majority of modern mathematicians, to be more sterile or more detached from all possible use, than those which concern speculation about the nature of numbers and research into their divisors. In this judgement the mathematicians of today differ greatly from the Ancients, who were accustomed to accord a great value to these speculations...For as well as it seeming to them that investigation of the truth was in itself laudable and worthy of human consciousness, they judged also, rightly, that by these researches the art of investigation could be extended, and that the faculties of the mind would become better able to deal with important questions. And in this opinion they were not deceived, for we have manifest proof of this in the considerable developments that have enriched Analysis since that epoch; in fact it appears entirely to be the case that science would never have achieved such a degree of perfection had the Ancients not put so much zeal into developing questions of this type, which the greater part of modern mathematicians despise so much on account of their sterility. - L. Euler (see [CS97])


## 1. Introduction and overview

For each natural number $n$, let $\sigma(n)$ be the sum of all the (positive) divisors of $n$, and let $s(n)$ be the sum of all the proper divisors. Here a proper divisor of $n$ is a divisor of $n$ other than $n$ itself, so that $s(n)=\sigma(n)-n$. The ancient Greeks partitioned the natural numbers according to whether $s(n)<n, s(n)=n$, or $s(n)>n$ (equivalently, $\sigma(n)<2 n, \sigma(n)=2 n$, or $\sigma(n)>2 n)$. Numbers $n$ of the first kind were termed deficient, numbers of the third kind abundant, and numbers of the second kind perfect.

Fast-forwarding to modern times, it is routine to verify by computer that among the first million natural numbers, 247,545 are abundant, 752,451 are deficient, and only 4 are perfect. This simple computation raises a number of questions: It seems that both the abundant and deficient numbers are relatively common. Do both of these sets constitute a well-defined proportion of the natural numbers? More precisely, is it true that the set of abundant numbers (or deficient numbers) possesses an asymptotic density? Given that we found just four perfect numbers up to $10^{6}$, should we expect infinitely many as we head out to infinity? The first four perfect numbers are

$$
6=2 \cdot 3, \quad 28=2^{2} \cdot 7, \quad 496=2^{4} \cdot 31, \quad \text { and } \quad 8128=2^{6} \cdot 127
$$

Are all perfect numbers even? Do they all only have two prime factors? The astute reader may have noticed that in our examples, all four factorizations have the form $2^{k}\left(2^{k+1}-1\right)$; does this continue?
1.1. Even perfect numbers. Let us turn to what is known about these questions. Suppose first that $2^{k+1}-1$ is a prime number. It was known already to Euclid that in this case the number $n:=2^{k}\left(2^{k+1}-1\right)$ is perfect, and this can be verified very quickly using the multiplicativity of the $\sigma$ function:

$$
\begin{aligned}
\sigma(n) & =\sigma\left(2^{k}\right) \sigma\left(2^{k+1}-1\right) \\
& =\left(1+2+4+\cdots+2^{k}\right)\left(1+\left(2^{k+1}-1\right)\right)=\left(2^{k+1}-1\right) 2^{k+1}=2 n,
\end{aligned}
$$

so that $s(n)=\sigma(n)-n=n$, i.e., $n$ is perfect.
Two thousand years later, Euler established a partial converse by proving that Euclid's rule accounts for every even perfect number. Here is a simple argument for this: Suppose that $n$ is an even perfect mumber and write $n=2^{k} q$, where $q$ is odd and $k \geq 1$. Then

$$
\begin{equation*}
2^{k+1} q=2 n=\sigma(n)=\sigma\left(2^{k}\right) \sigma(q)=\left(2^{k+1}-1\right) \sigma(q) . \tag{8.1}
\end{equation*}
$$

Because $2^{k+1}-1$ and $2^{k+1}$ are coprime, it must be that $\left(2^{k+1}-1\right) \mid q$, so that we can write $q=\left(2^{k+1}-1\right) r$. Substituting this expression for $q$ into
(8.1), we obtain (upon canceling $2^{k+1}-1$ from both sides) that

$$
\begin{equation*}
2^{k+1} r=\sigma(q) . \tag{8.2}
\end{equation*}
$$

This forces us to have $r=1$, since otherwise $1, r$, and $\left(2^{k+1}-1\right) r$ are distinct divisors of $q$ which sum to more than

$$
\left(2^{k+1}-1\right) r+r=2^{k+1} r=\sigma(q) .
$$

Hence $q=2^{k+1}-1$. Moreover, putting $r=1$ in (8.2), we obtain

$$
\sigma(q)=2^{k+1} .
$$

So $\sigma(q)=q+1$. But this implies that $q$ is prime. So $n=2^{k}\left(2^{k+1}-1\right)$, where the second factor is prime, and this is exactly what we set out to show.

Summarizing, we have proved the following classical result:
Theorem 8.1 (Euclid-Euler). If $2^{k+1}-1$ is prime, then $2^{k}\left(2^{k+1}-1\right)$ is a perfect number. Conversely, if $n$ is an even perfect number, then $n=$ $2^{k}\left(2^{k+1}-1\right)$ for some $k \geq 1$ for which $2^{k+1}-1$ is prime.

The Euclid-Euler classification more or less closes the book on even perfect numbers. Of course it does not single-handedly answer all of the many questions one might have about these numbers, but it shows that such questions may be thought of as questions about primes of the form $2^{k+1}-1$ (so-called Mersenne primes). These new questions may in turn prove intractable, but the blame now rests with the analytic number theorists and not the investigator of perfect numbers. As an example of this process of translation, consider the question of how many even perfect numbers there are up to $x$. In Chapter 3, we suggested (Conjecture 3.20) that $2^{m}-1$ is prime for $(1+o(1)) e^{\gamma} \log x / \log 2$ values of $m \leq x$. So from the Euclid-Euler result, we find that the number of even perfect numbers up to $x$ should be

$$
\sim \frac{e^{\gamma}}{\log 2} \log \log x
$$

1.2. Odd perfect numbers. So what about odd perfect numbers? Here much less is known; in particular, not a single example has ever been discovered. One of the earliest results of substance is due to Euler, who showed that the factorization of a hypothetical odd perfect number must take a certain peculiar form: Suppose that $n$ is an odd perfect number, and write the prime factorization of $n$ in the form $n=\prod_{i=0}^{k} p_{i}^{f_{i}}$. Since $n$ is odd,

$$
2 \| 2 n=\sigma(n)=\prod_{i=0}^{k} \sigma\left(p_{i}^{f_{i}}\right) .
$$

As a consequence, each term $\sigma\left(p_{i}^{f_{i}}\right)$ in the product is odd except for a single exceptional value of $i$, where $2 \| \sigma\left(p_{i}^{f_{i}}\right)$. By relabeling if necessary, we can
assume $i=0$ corresponds to the special term. Since each of the primes $p_{i}$ is odd, we have

$$
\sigma\left(p_{i}^{f_{i}}\right)=1+p_{i}+p_{i}^{2}+\cdots+p_{i}^{f_{i}} \equiv f_{i}+1(\bmod 2),
$$

and so $f_{i}$ must be even for every $1 \leq i \leq k$. For $i=0$, the condition $2 \| \sigma\left(p_{i}^{f_{i}}\right)$ says that $\sigma\left(p_{0}^{f_{0}}\right) \equiv 2(\bmod 4)$. But it is easy to check that this happens only when $p_{0} \equiv f_{0} \equiv 1(\bmod 4)$. We have thus proved (writing $e=f_{0}$ and $e_{i}=\frac{1}{2} f_{i}$ for $\left.1 \leq i \leq k\right)$ :

Theorem 8.2 (Euler). Every odd perfect number has the form $p^{e} \prod_{i=1}^{k} p_{i}^{2 e_{i}}$, where $p$ and the $p_{i}$ are distinct primes, and $p \equiv e \equiv 1(\bmod 4)$.

Since the time of Euler, several mathematicians have obtained other results on what an odd perfect number must look like if one exists. Here are four results representative of the current state-of-the-art: If $n$ is an odd perfect number, then:

- $n$ has more than 300 decimal digits (Brent, Cohen \& te Riele [BCtR91]),
- $n$ has a prime factor larger than $10^{8}$ (Goto \& Ohno [GO08]),
- $n$ has at least 9 distinct prime factors (Nielsen [Nie07]),
- $n$ has at least 75 prime factors, counted with multiplicity (Hare [Har07]).
While at their core the arguments of these four papers are elementary, in each case the proofs require extensive computer work. We will not prove these results here. Instead we focus our discussion of odd perfect numbers on two theorems not about the structure of individual odd perfect numbers, but about the set of odd perfect numbers as a whole. The first is the following "finiteness theorem" due to Dickson [Dic13a]:

Theorem 8.3. For each fixed $k \in \mathbf{N}$, there are only finitely many odd perfect numbers with precisely $k$ distinct prime factors.

Theorem 8.3 shows that odd perfect numbers behave quite differently from even perfect numbers, where each (of the probably infinitely many examples) has exactly two distinct prime factors.

Upon reading the statement of Theorem 8.3, it is natural to think that the result of Nielsen quoted above has been reduced to a finite check. But this is not the case: The proof we will give of Theorem 8.3 in $\S 8.3$ is ineffective, in that while it shows that there are at most finitely many examples for each fixed value of $k$, it does not yield any finite procedure for finding all of them. Doing a bit more work, one can prove an effective version of

Theorem 8.3. Indeed, Pomerance [Pom77a] has shown that an odd perfect number with $k$ distinct prime factors is necessarily less than

$$
(4 k)^{(4 k)^{2^{k^{2}}}}
$$

so that (in principle) one can simply test all candidates up to this bound! Heath-Brown [HB94] has shown that this gargantuan bound can be replaced with the (still astronomical) $4^{4^{k}}$, and Nielsen [Nie03] has reduced this further to $2^{4^{k}}$. (Of course, since $2^{4^{8}}>2 \cdot 10^{19728}$, this is not how Nielsen shows that an odd perfect number has at least 9 prime factors; more cunning is required!)

Our second theorem addresses the distribution of odd perfect numbers. In Theorem 6.4, we saw already how to deduce from Theorem 8.2 that the set of odd perfect numbers has density zero. That is a rather weak result, and it is easy to do substantially better: From Theorem 8.2, every odd perfect number $n$ has the form $p^{e} m^{2}$, where $\operatorname{gcd}(p, m)=1$. If $n \leq x$, then clearly $m \leq \sqrt{x}$. So let us fix a natural number $m \leq \sqrt{x}$ and ask for a prime power $p^{e}$ with $\operatorname{gcd}(p, m)=1$ for which $p^{e} m^{2}$ is perfect. In that case,

$$
\sigma\left(p^{e} m^{2}\right)=2 p^{e} m^{2}, \quad \text { so that } \quad \frac{\sigma\left(p^{e}\right)}{p^{e}}=\frac{2 m^{2}}{\sigma\left(m^{2}\right)}
$$

But as $p^{e}$ ranges over prime powers, the numbers $\sigma\left(p^{e}\right) / p^{e}$ are all distinct; the simplest way to see this is to observe that $\sigma\left(p^{e}\right) / p^{e}$ is already a fraction in lowest terms. So there can be at most one prime power $p^{e}$ (with $p \nmid m$ ) making $p^{e} m^{2}$ perfect, and we obtain immediately that there are at most $x^{1 / 2}$ odd perfect numbers $n \leq x$. This simple argument is due to Hornfeck [Hor55]. Later, in joint work with Wirsing, Hornfeck established [HW57] that the number of odd perfect numbers up to $x$ is $O_{\epsilon}\left(x^{\epsilon}\right)$ for each $\epsilon>0$. The strongest known result in this direction is due to Wirsing [Wir59]:

Theorem 8.4. There is an absolute constant $W>0$ with the property that the number of perfect numbers $n \leq x$ is smaller than $x^{W / \log \log x}$ for every $x \geq 3$.

We will give Wirsing's proof in $\S 3$. In that section we also include a heuristic argument, due to Pomerance, suggesting that probably there aren't any odd perfect numbers at all.
1.3. The density of the abundant numbers. So far we have yet to answer the very first question we posed: Does the set of abundant numbers have an asymptotic density? The answer to this question is "yes", and in fact much more is true. The following beautiful result is due to Davenport [Dav33]; we give an elementary proof (essentially due to Erdős) in $\S 4$.

Theorem 8.5. For each real number $u$, the set of natural numbers $n$ for which $\sigma(n) / n \leq u$ possesses an asymptotic density. Calling this density $D(u)$, the function $D$ is continuous on all of $\mathbf{R}$ and satisfies $D(1)=0$ and $\lim _{u \rightarrow \infty} D(u)=1$.

The function $D(u)$ is known as the distribution function for $\sigma(n) / n$.
Since (as discussed above) the set of perfect numbers has density zero, ${ }^{1}$ it is immediate from Theorem 8.5 that the deficient numbers have density $D(2)$ and the abundant numbers have density $1-D(2)$. M. Kobayashi, improving earlier results of Behrend [Beh33], Wall et al. [Wal72, WCJ72], and Deléglise [Del98], shows in his Ph.D. thesis [Kob10] that

$$
0.24761<1-D(2)<0.24767
$$

So just under 1 in 4 natural numbers are abundant. Precise numerical values of $D(2)$ and $1-D(2)$ are not important for the rest of this chapter, but it will be useful to keep in mind that the abundant numbers have positive density. (This is obvious once one knows that the density exists, since, e.g., it is easily shown that every multiple of 12 is abundant.)
1.4. Aliquot sequences and sociable numbers. In the remainder of this chapter we broaden our study to include certain relatives of the perfect numbers. Say that two (distinct) natural numbers $m$ and $n$ form an amicable pair if each is the sum of the proper divisors of the other, i.e., if $s(m)=n$ and $s(n)=m$. In this case both $m$ and $n$ are called amicable. For example, 220 and 284 form an amicable pair, since

$$
\begin{aligned}
& s(284)=1+2+4+71+142=220, \quad \text { while } \\
& \quad s(220)=1+2+4+5+10+11+20+22+44+55+110=284 .
\end{aligned}
$$

The study of amicable numbers goes back to the Pythagoreans, but still many of the simplest questions remain unanswered. For example, while there are over 12 million examples of amicable pairs known ([Ped]; see also [GPtR04]), we have no proof that there are infinitely many.

To understand the relation between amicable numbers and perfect numbers, it is illuminating to bring into play the concept of an aliquot sequence. Let $s_{k}$ be the $k$ th iterate of $s$, defined as follows: $s_{0}(n)=n$, and if $k \geq 0$ and $s_{k}(n)>0$, then $s_{k+1}(n):=s\left(s_{k}(n)\right)$. The sequence of iterates $n, s(n), s_{2}(n), \ldots$ is called the aliquot sequence at $n$. For example, if $n=24$, we obtain $24,36,55,17,1,0$, and so the sequence terminates. However, if $n=25$, the sequence is $25,6,6,6, \ldots$, so is eventually periodic. A conjecture of Catalan [Cat88] (as corrected by Dickson [Dic13b]) asserts that if

[^4]$n$ is any natural number, then the aliquot sequence at $n$ is always either terminating or eventually periodic.

The Catalan-Dickson conjecture has been verified by computer to hold for all $n<276$. But when $n=276$, the corresponding aliquot sequence has been computed to well over a thousand terms without any repetition. Guy \& Selfridge [GS75] have suggested that the early initial evidence for the Catalan-Dickson conjecture is deceptive, and that infinitely many aliquot sequences, perhaps most of those that start at an even value of $n$, tend to infinity. We will not enter into this controversy here.

We say that the natural number $n$ is sociable if the aliquot sequence at $n$ is purely periodic. In this case we call the length $k$ of the period the order of $n$, and the set $\left\{n, s(n), s_{2}(n), \ldots, s_{k-1}(n)\right\}$ is called a sociable cycle of length (or order) $k$. For example, 7169104 starts a sociable cycle of length 4 , since under repeated application of $s$,

$$
7169104 \mapsto 7538660 \mapsto 8292568 \mapsto 7520432 \mapsto 7169104 \mapsto \cdots
$$

At present, there are 175 known examples of sociable cycles of length $>2$ [Moe]; of these, all but 10 have length 4.

It is reasonable to wonder what can be said, in general, about the distribution of sociable numbers. The following theorem is due to Erdős:

Theorem 8.6. For each fixed $k \in \mathbf{N}$, the set of sociable numbers of order $k$ has asymptotic density zero.

Probably much more than Theorem 8.6 is true; the authors of [KPP09] conjecture that the set of all sociable numbers has density zero and prove that this holds if we discard the odd abundant members of this set.

It should be noted that when $k=2$, Pomerance [Pom81] has proved a much stronger upper bound than that furnished by Theorem 8.6:
$\star$ Theorem 8.7. The number of amicable numbers $n \leq x$ is smaller than $x / \exp \left((\log x)^{1 / 3}\right)$ for all sufficiently large $x$.

No result of comparable strength is known when $k>2$.

## 2. Proof of Dickson's finiteness theorem

Lemma 8.8. Let $k \in \mathbf{N}$. Suppose that $A$ is an infinite, strictly increasing sequence of natural numbers each of which has precisely $k$ distinct prime divisors. Then we may extract from $A$ an infinite subsequence $\left\{n_{j}\right\}_{j=1}^{\infty}$, where each $n_{j}$ has the form

$$
\begin{equation*}
n_{j}:=p_{1}^{e_{1}} \cdots p_{r}^{e_{r}} p_{r+1}^{e_{r+1, j}} \cdots p_{t}^{e_{t, j}} p_{t+1, j}^{e_{t+1, j}} \cdots p_{k, j}^{e_{k, j}} \tag{8.3}
\end{equation*}
$$

and where
(i) $p_{i}^{e_{i}}$ is fixed independently of $j$ for $1 \leq i \leq r$,
(ii) $p_{i}$ is fixed independently of $j$ and $e_{i, j} \rightarrow \infty$ as $j \rightarrow \infty$, for each $r<i \leq t$, and
(iii) $p_{i, j} \rightarrow \infty$ as $j \rightarrow \infty$, for $t<i \leq k$.

Proof. If there is an infinite subsequence of $A$ all of whose terms are exactly divisible by a fixed prime power $p_{1}^{e_{1}}$, pass to this subsequence. If there is an infinite subsequence of remaining terms exactly divisible by some other prime power $p_{2}^{e_{2}}$, then pass to this subsequence. Continuing, we eventually arrive at an infinite sequence all of whose terms are exactly divisible by $p_{1}^{e_{1}}, \ldots, p_{r}^{e_{r}}$ (say), and which does not have any infinite subsequence of integers whose canonical factorizations contain a fixed prime power other than $p_{1}^{e_{1}}, \ldots, p_{r}^{e_{r}}$. (This process necessarily terminates in $r \leq k$ steps. Of course it is also possible that it never starts, i.e., that $r=0$.)

If at this point our sequence has an infinite subsequence all of whose terms are divisible by a fixed prime $p_{r+1}$ different from $p_{1}, \ldots, p_{r}$, then pass to this subsequence. Note that the exponent of $p_{r+1}$ along the terms of this subsequence must tend to infinity to avoid contradicting the conclusion of the last paragraph. If our sequence has an infinite subsequence all of whose terms are divisible by the fixed prime $p_{r+2} \notin\left\{p_{1}, \ldots, p_{r+1}\right\}$, pass to this subsequence. Continue this process as long as possible, ending with (say) $p_{t}$. Then our final sequence has all of the properties specified in Lemma 8.8.

Lemma 8.9. For every natural number $n$, we have $\sigma(n) / n=\sigma_{-1}(n)$, where

$$
\sigma_{-1}(n):=\sum_{d \mid n} \frac{1}{d} .
$$

Consequently, if $m$ and $n$ are two natural numbers for which $m \mid n$, then $\sigma(m) / m \leq \sigma(n) / n$ with equality only if $m=n$.

Proof. We have $\sigma(n) / n=(1 / n) \sum_{d \mid n} d=\sum_{d \mid n}(n / d)^{-1}=\sigma_{-1}(n)$, since $n / d$ runs over all the divisors of $n$ as $d$ does. The rest of the lemma is now obvious.

Lemma 8.9 implies, in particular, that it is impossible for one perfect number to properly divide another.

Proof of Theorem 8.3 (Shapiro [Sha49b]). Suppose that there are infinitely many odd perfect numbers with exactly $k$ distinct prime factors, and let $A$ be the sequence of such numbers in increasing order. Use Lemma 8.8 to extract an infinite subsequence $n_{1}<n_{2}<n_{3}<\cdots$ of $A$ whose factorizations have the form (8.3). Applying $\sigma_{-1}$ to both sides of (8.3), we find that
for each $j=1,2,3, \ldots$,

$$
\begin{aligned}
2=\sigma_{-1}\left(n_{j}\right) & =\prod_{i=1}^{r} \sigma_{-1}\left(p_{i}^{e_{i}}\right) \prod_{i=r+1}^{t} \sigma_{-1}\left(p_{i}^{e_{i, j}}\right) \prod_{i=t+1}^{k} \sigma_{-1}\left(p_{i, j}^{e_{i, j}}\right) \\
& =\prod_{i=1}^{r} \frac{p_{i}^{e_{i}+1}-1}{p_{i}^{e_{i}}\left(p_{i}-1\right)} \prod_{i=r+1}^{t}\left(1+\frac{1}{p_{i}}+\cdots+\frac{1}{p_{i}^{e_{i, j}}}\right) \prod_{i=t+1}^{k} \sigma_{-1}\left(p_{i, j}^{e_{i, j}}\right) .
\end{aligned}
$$

Letting $j \rightarrow \infty$, we find (referring back to the statement of Lemma 8.8) that

$$
2=\prod_{i=1}^{r} \frac{p_{i}^{e_{i}+1}-1}{p_{i}^{e_{i}}\left(p_{i}-1\right)} \prod_{i=r+1}^{t} \frac{p_{i}}{p_{i}-1},
$$

so that

$$
\begin{equation*}
2 \prod_{i=1}^{r} p_{i}^{e_{i}} \prod_{i=r+1}^{t}\left(p_{i}-1\right)=\prod_{i=1}^{r} \frac{p_{i}^{e_{i}+1}-1}{p_{i}-1} \prod_{i=r+1}^{t} p_{i} . \tag{8.4}
\end{equation*}
$$

Since $p_{r+1}, p_{r+2}, \ldots, p_{t}$ are odd primes distinct from $p_{1}, \ldots, p_{r}$, (8.4) implies that $\prod_{i=r+1}^{t} p_{i} \mid \prod_{i=r+1}^{t}\left(p_{i}-1\right)$. This is only possible if both products are empty, i.e., if $r=t$. In this case,

$$
2 \prod_{i=1}^{r} p_{i}^{e_{i}}=\prod_{i=1}^{r} \frac{p_{i}^{e_{i}+1}-1}{p_{i}-1}
$$

which says that $n:=\prod_{i=1}^{r} p_{i}^{e_{i}}$ is perfect; but this is impossible, since $n$ divides every $n_{j}$ and no perfect number can properly divide another.

## 3. How rare are odd perfect numbers?

3.1. Proof of Wirsing's theorem. We need two combinatorial lemmas before we can prove Theorem 8.4:

Lemma 8.10. Let $M$ be a nonnegative integer. Then there are exactly $2^{M}$ solutions to the inequality

$$
e_{1}+e_{2}+\cdots+e_{k} \leq M,
$$

where $k \geq 0$ and the $e_{i}$ are positive integers. Here the empty sum is counted as a solution corresponding to $k=0$.

Proof. Define the formal power series $P(T)$ by putting $P(T):=T+T^{2}+$ $T^{3}+\cdots$. Then $P(T)=T /(1-T)$. Moreover, the number of solutions in positive integers $e_{1}, \ldots, e_{k}$ to $e_{1}+\cdots+e_{k}=m$ is given by the coefficient of $T^{m}$ in

$$
1+P(T)+P(T)^{2}+P(T)^{3}+\cdots=\frac{1}{1-P(T)}=\frac{1-T}{1-2 T} .
$$

Consequently, the quantity described in the lemma statement is given by the coefficient of $T^{M}$ in

$$
\left(1+T+T^{2}+\ldots\right) \frac{1-T}{1-2 T}=\frac{1}{1-T} \frac{1-T}{1-2 T}=\frac{1}{1-2 T},
$$

which is just $2^{M}$, as claimed.
Lemma 8.11. Let $M$ and $k$ be nonnegative integers. Then the inequality

$$
e_{1}+e_{2}+\cdots+e_{k} \leq M
$$

has exactly $\binom{M+k}{M} \leq 2^{M+k}$ solutions in nonnegative integers $e_{1}, e_{2}, \ldots, e_{k}$.
Proof. The number of solutions to the inequality of the lemma is the same as the number of solutions in nonnegative $e_{i}$ to the equation $e_{0}+e_{1}+\cdots+e_{k}=$ $M$. This is given by the coefficient of $T^{M}$ in the power series

$$
\left(1+T+T^{2}+T^{3}+\cdots\right)^{k+1}=(1-T)^{-(k+1)},
$$

which by the binomial theorem is precisely

$$
(-1)^{M}\binom{-k-1}{M}=\binom{M+k}{M}
$$

as claimed. The upper bound $\binom{M+k}{M} \leq 2^{M+k}$ is obvious, since $\binom{M+k}{M}$ is a summand in the binomial expansion of $(1+1)^{M+k}$.

Proof of Theorem 8.4. For each perfect number $n \leq x$, we write $n=A Q$, where

$$
\begin{equation*}
A:=\prod_{\substack{p^{e_{p}} \| n \\ p>\log x}} p^{e_{p}} \quad \text { and } \quad Q:=\prod_{\substack{p^{e_{p} \| n} \\ p \leq \log x}} p^{e_{p}} . \tag{8.5}
\end{equation*}
$$

Thus $Q$ represents the $(\log x)$-smooth part of $n .{ }^{2}$ Loosely speaking, we will show that $Q$ essentially determines $A$, and so also essentially determines $n$. Theorem 8.4 will then follow from an upper bound on the number of $(\log x)$-smooth integers $Q \leq x$.

Fix a $(\log x)$-smooth integer $Q \leq x$. Let us suppose $n \leq x$ is perfect, where $n=A Q$ and every prime factor of $A$ exceeds $\log x$. Then $A$ can have at most $\log x / \log \log x$ distinct prime factors, and so

$$
\begin{aligned}
\frac{A}{\sigma(A)} & \geq \prod_{\substack{\log x<p \leq x \\
p \mid A}}\left(1-\frac{1}{p}\right) \geq 1-\sum_{\substack{\log x<p \leq x \\
p \mid A}} \frac{1}{p} \\
& \geq 1-\frac{1}{\log x} \frac{\log x}{\log \log x}=1-\frac{1}{\log \log x}>1 / 2
\end{aligned}
$$

[^5]if $x$ is large (which we assume). Since
\[

$$
\begin{equation*}
\sigma(A) \sigma(Q)=\sigma(n)=2 n=2 A Q \tag{8.6}
\end{equation*}
$$

\]

we have

$$
\begin{equation*}
Q<\frac{2 A}{\sigma(A)} Q=\sigma(Q) \leq 2 Q \tag{8.7}
\end{equation*}
$$

with equality on the right only if $A=1$. Thus if $A \neq 1$, then $\sigma(Q) \nmid 2 Q$, so that there is a prime dividing $\sigma(Q)$ to a higher power than it divides $2 Q$. Let $p_{1}$ be the least such prime. It follows from (8.6) that $p_{1}^{e_{1}} \| A$ for a certain exponent $e_{1} \geq 1$. Now write

$$
n=A^{\prime} Q^{\prime}, \quad \text { where } \quad A^{\prime}:=\frac{A}{p_{1}^{e_{1}}}, \quad Q^{\prime}:=Q p_{1}^{e_{1}}
$$

Then $A^{\prime} / \sigma\left(A^{\prime}\right) \geq A / \sigma(A)>1 / 2$ and both (8.6) and (8.7) hold with $A$ and $Q$ replaced by $A^{\prime}$ and $Q^{\prime}$ (respectively). Repeating the above argument, we find that if $A^{\prime} \neq 1$, then there exists a prime dividing $\sigma\left(Q^{\prime}\right)$ to a higher power than it divides $2 Q^{\prime}$. Letting $p_{2}$ be the smallest such prime, we have that $p_{2}^{e_{2}} \| A^{\prime}$ for a certain exponent $e_{2} \geq 1$. We then set $A^{\prime \prime}:=A^{\prime} / p_{2}^{e_{2}}, Q^{\prime \prime}:=$ $Q^{\prime} p_{2}^{e_{2}}$, and continue. This process eventually terminates and we obtain a factorization of the form

$$
A=p_{1}^{e_{1}} p_{2}^{e_{2}} \cdots p_{t}^{e_{t}} .
$$

Notice that the prime $p_{1}$ is completely determined by $Q$, while for $i>1$, the prime $p_{i}$ is completely determined by $Q$ and $e_{1}, e_{2}, \ldots, e_{i-1}$. So, for fixed $Q$, the number $A$ is completely determined by the sequence of exponents $e_{1}, \ldots, e_{t}$. Since each prime divisor of $A$ exceeds $\log x$, we have

$$
e_{1}+\cdots+e_{t} \leq \frac{\log x}{\log \log x}
$$

It now follows from Lemma 8.10 that for each $Q$, there are at most

$$
\begin{equation*}
2^{\log x / \log \log x}=x^{\log 2 / \log \log x} \tag{8.8}
\end{equation*}
$$

choices for $A$.
It remains to estimate the number $(\log x)$-smooth natural numbers $Q \leq$ $x$. For each such $Q$, put $Q=Q_{1} Q_{2}$, where

$$
Q_{1}:=\prod_{\substack{p^{e_{p}} \| Q \\ \sqrt{\log x}<p \leq \log x}} p^{e_{p}} \quad \text { and } \quad Q_{2}:=\prod_{\substack{p^{e_{p}} \| Q \\ p \leq \sqrt{\log x}}} p^{e_{p}}
$$

Because $Q_{1} \leq Q \leq x$, the exponents $e_{p}$ appearing in the factorization of $Q_{1}$ must satisfy

$$
\sum_{\sqrt{\log x}<p \leq \log x} e_{p} \leq \frac{\log x}{\log \sqrt{\log x}}=2 \frac{\log x}{\log \log x}
$$

The number of summands here is bounded by $\pi(\log x)$, which by the results of Chebyshev is at most $K \log x / \log \log x$ for a certain constant $K>0$. So by Lemma 8.11, the number of possibilities for $Q_{1}$ is at most

$$
\begin{equation*}
2^{(K+2) \log x / \log \log x}=x^{(K+2) \log 2 / \log \log x} . \tag{8.9}
\end{equation*}
$$

Since also $Q_{2} \leq x$, the exponent $e_{p}$ of each prime appearing in the definition of $Q_{2}$ is trivially $\leq \log x / \log 2$. Thus the number of possibilities for $Q_{2}$ is bounded by

$$
\begin{align*}
\prod_{p \leq \sqrt{\log x}}\left(1+\frac{\log x}{\log 2}\right) & \leq\left(1+\frac{\log x}{\log 2}\right)^{\sqrt{\log x}} \\
& \leq \exp (2 \sqrt{\log x} \log \log x)=\exp (o(\log x / \log \log x)) \tag{8.10}
\end{align*}
$$

From (8.9) and (8.10), the total number of $(\log x)$-smooth $Q \leq x$ is at most $x^{((K+2) \log 2+o(1)) / \log \log x}$. So from (8.8), if $W>(K+3) \log 2$, then the number of perfect numbers $\leq x$ is at most most $x^{W / \log \log x}$ for all sufficiently large values of $x$. Adjusting the value of $W$ if necessary, this can be made to hold for all $x \geq 3$.
3.2. A heuristic. Theorem 8.2 tells us that every odd perfect number can be written in the form $p^{e} m^{2}$, where $p \equiv e \equiv 1(\bmod 4)$ and $\operatorname{gcd}(p, m)=1$. Call a number $n$ of this form a candidate, and say that $n$ is successful if $n$ is actually an odd perfect number, i.e., if

$$
2 p^{e} m^{2}=2 n=\sigma(n)=\sigma\left(p^{e}\right) \sigma\left(m^{2}\right) .
$$

Let us attempt to estimate the odds that a given $m$ corresponds to a successful candidate $n=p^{e} m^{2}$. (Strictly speaking this is nonsense, since such an $m$ either does or doesn't correspond to such an $n$; there are no "odds" about it. But it is a useful bit of nonsense!) Since $\operatorname{gcd}\left(p^{e}, \sigma\left(p^{e}\right)\right)=1$, if $n$ is successful, then $p^{e} \| \sigma\left(m^{2}\right)$. The number of exact prime power divisors of $\sigma\left(m^{2}\right)$ is trivially at most $\log \sigma\left(m^{2}\right) / \log 2$. Since

$$
\sigma\left(m^{2}\right) \leq \sum_{d \leq m^{2}} d \leq m^{2} \cdot m^{2}=m^{4}
$$

there are at most $4 \log m / \log 2$ possibilities for $p^{e}$. Supposing now that $p^{e}$ does exactly divide $m^{2}$, for $n=p^{e} m^{2}$ to be successful, we also need that

$$
m^{2} \left\lvert\, \sigma\left(p^{e}\right) \frac{\sigma\left(m^{2}\right)}{p^{e}}\right.,
$$

which we might expect to hold with "probability" $1 / m^{2}$.
The upshot is that for a given value of $m$, the "probability" that there is a successful candidate of the form $p^{e} m^{2}$ is at most $(4 \log 2) \log m / m^{2}$. Since
the sum $\sum_{m \geq 1}(4 \log 2) \log m / m^{2}$ converges, we expect that there are only finitely many successful candidates (odd perfect numbers).

We can take this a bit further. Suppose $n$ is an odd perfect number, and write $n=p^{e} m^{2}$ as above. Since $m^{2}$ is a proper divisor of $n$, we have $\sigma\left(m^{2}\right) / m^{2}<\sigma(n) / n=2$. Since also $p^{e}$ divides $\sigma\left(m^{2}\right)$, it follows that

$$
2 m^{4}>\sigma\left(m^{2}\right) m^{2} \geq p^{e} m^{2}=n>10^{300}
$$

using the result of Brent, Cohen, and te Riele mentioned on p. 250. Thus $m>2^{-1 / 4} \cdot 10^{75}$. If we compute $\sum(4 \log 2) \log m / m^{2}$ over these values of $m$, we obtain an upper bound of less than $10^{-70}$ for the expected total count of odd perfect numbers. So it seems highly unlikely that any example exists.

This is (a slight variant of) an unpublished argument of Pomerance.

## 4. The distribution function of $\sigma(n) / n$

Theorem 8.5 asserts that for each real $u$, the density of the set of $n$ with $\sigma(n) / n \leq u$ exists; moreover, calling this density $D(u)$, we have that $D(u)$ is a continuous function of $u, D(1)=0$, and $\lim _{u \rightarrow \infty} D(u)=1$. Owing to Lemma 8.9, we may replace " $\sigma(n) / n$ " in this statement with " $\sigma_{-1}(n)$ ", which will prove convenient both notationally and psychologically.

For each $B>0$, we define the arithmetic function $\sigma_{-1}^{B}$ by putting

$$
\sigma_{-1}^{B}(n):=\sum_{\substack{d|n \\ p| d \Rightarrow p \leq B}} \frac{1}{d} .
$$

In other words, $\sigma_{-1}^{B}$ is obtained by restricting the sum defining $\sigma_{-1}$ to $B$ smooth divisors of $n$. We also set $F^{B}(n)$ equal to the $B$-smooth part of $n$, i.e., $F^{B}(n):=\prod_{p^{e} \| n, p \leq B} p^{e}$. Note that with these definitions, we have $\sigma_{-1}^{B}(n)=\sigma_{-1}\left(F^{B}(n)\right)$. Define
$\mathcal{N}(x, u):=\left\{n \leq x: \sigma_{-1}(n) \leq u\right\} \quad$ and $\quad \mathcal{N}^{B}(x, u):=\left\{n \leq x: \sigma_{-1}^{B}(n) \leq u\right\}$, and set $N(x, u):=\# \mathcal{N}(x, u)$ and $N^{B}(x, u):=\# \mathcal{N}^{B}(x, u)$.

We begin the proof of Theorem 8.5 by demonstrating a partial analogue of that result for the functions $\sigma_{-1}^{B}$ :

Lemma 8.12. Let $B>0$. For each real $u$, the quantity $N^{B}(x, u) / x$ tends to a limit, say $D^{B}(u)$, as $x \rightarrow \infty$.

Proof. Let $S$ be the collection of $B$-smooth numbers $m$ with $\sigma_{-1}(m) \leq u$. For a natural number $n$, we have $\sigma_{-1}^{B}(n) \leq u$ precisely when $F^{B}(n)=m$ for some $m \in S$. For each $m \in S$, the set of natural numbers $n$ with $F^{B}(n)=m$ possesses an asymptotic density, since this set is just the union of certain
residue classes modulo $m \prod_{p \leq B} p$. Denote this density by $d_{m}$. We claim that $N^{B}(x, u) / x \rightarrow \sum_{m \in S} d_{m}$ as $x \rightarrow \infty$.

For the proof, let $z$ be a positive real parameter. Since $\sigma_{-1}^{B}(n) \leq u$ whenever $F^{B}(n) \in S \cap[1, z]$, it is clear that

$$
\begin{equation*}
\liminf _{x \rightarrow \infty} \frac{N^{B}(x, u)}{x} \geq \sum_{\substack{m \in S \\ m \leq z}} d_{m} \tag{8.11}
\end{equation*}
$$

On the other hand, if $\sigma_{-1}^{B}(n) \leq u$, then either $F^{B}(n) \in S \cap[1, z]$, or $n$ is divisible by some $m \in S$ with $m>z$. So

$$
\begin{equation*}
\limsup _{x \rightarrow \infty} \frac{N^{B}(x, u)}{x} \leq \sum_{\substack{m \in S \\ m \leq z}} d_{m}+\sum_{\substack{m \in S \\ m>z}} \frac{1}{m} \tag{8.12}
\end{equation*}
$$

Since $\sum_{m \in S} m^{-1} \leq \sum_{m B \text {-smooth }} m^{-1}=\prod_{p \leq B}(1-1 / p)^{-1}<\infty$, the final sum in (8.12) is the tail of a convergent series. So the desired equality $D^{B}(u)=\sum_{m \in S} d_{m}$ follows by letting $z \rightarrow \infty$ in (8.11) and (8.12).

Lemma 8.13. Let $\mathcal{P}$ be a set of primes for which $\sum_{p \in \mathcal{P}} p^{-1}$ diverges. For each $\epsilon>0$, there is a $z>0$ for which the following holds: For all $n$ outside of a set of density $<\epsilon$, there is a prime $p \in \mathcal{P} \cap[2, z]$ for which $p \| n$.

Proof. The relation $p \| n$ holds precisely when $n$ falls into one of the $p-1$ residue classes $p, 2 p, \ldots,(p-1) p \bmod p^{2}$. So by the Chinese remainder theorem and the principle of inclusion-exclusion, the set of $n$ exactly divisible by none of the primes $p \in \mathcal{P} \cap[2, z]$ has density

$$
\prod_{\substack{p \in \mathcal{P} \\ p \leq z}}\left(1-\frac{p-1}{p^{2}}\right)<3 \exp \left(-\sum_{\substack{p \in \mathcal{P} \\ p \leq z}} \frac{1}{p}\right)
$$

which for large values of $z$ is less than $\epsilon$.
Lemma 8.14. Let $u$ be any real number. As $\delta \downarrow 0$, the upper density of the set of $n$ with $u-\delta<\sigma_{-1}(n)<u+\delta$ tends to zero.

Proof. Since the image of $\sigma_{-1}$ is contained in $[1, \infty)$, we may assume that $u \geq 1$. Let $\epsilon>0$, and fix a real number $B>0$ with $1 / B<\epsilon$. By Lemma 8.13, we can fix $z$ so that if $p_{1}<p_{2}<\cdots<p_{k}$ is the list of primes in the interval $(B, z]$, then all $n$ outside of an exceptional set of density $<\epsilon$ are exactly divisible by at least one of $p_{1}, \ldots, p_{k}$.

Let

$$
\mathcal{N}(x):=\left\{n \leq x: u-\delta<\sigma_{-1}(n)<u+\delta\right\} .
$$

For each $n \in \mathcal{N}(x)$ not in the exceptional set described above, fix a prime $p_{i}$ (with $1 \leq i \leq k$ ) exactly dividing $n$ and form the quotient $n / p_{i}$. We claim that if $\delta>0$ is chosen sufficiently small depending on $\epsilon$, then all of the quotients $n / p_{i}$ are distinct. Since each such quotient is at most $x / B$, for large $x$ this implies

$$
\# \mathcal{N}(x)<\epsilon x+x / B<2 \epsilon x
$$

which proves the lemma.
To establish the claim, suppose that $n$ and $n^{\prime}$ are distinct elements of $\mathcal{N}(x)$, that $p_{i} \| n$ and $p_{j} \| n^{\prime}$ (where $1 \leq i, j \leq k$ ), and that $n / p_{i}=n^{\prime} / p_{j}$. Clearly $i \neq j$. Moreover,

$$
\frac{\sigma_{-1}(n)}{\sigma_{-1}\left(p_{i}\right)}=\sigma_{-1}\left(n / p_{i}\right)=\sigma_{-1}\left(n^{\prime} / p_{j}\right)=\frac{\sigma_{-1}\left(n^{\prime}\right)}{\sigma_{-1}\left(p_{j}\right)},
$$

which implies that

$$
\frac{u-\delta}{u+\delta} \leq \frac{\sigma_{-1}(n)}{\sigma_{-1}\left(n^{\prime}\right)}=\frac{\sigma_{-1}\left(p_{i}\right)}{\sigma_{-1}\left(p_{j}\right)} \leq \frac{u+\delta}{u-\delta} .
$$

Thus, assuming $\delta<1 / 2$,

$$
\left|\frac{\sigma_{-1}\left(p_{i}\right)}{\sigma_{-1}\left(p_{j}\right)}-1\right| \leq \frac{2 \delta}{u-\delta}<4 \delta .
$$

(Recall that $u \geq 1$.) But this is impossible for sufficiently small values of $\delta$, since the numbers $\sigma_{-1}\left(p_{1}\right), \ldots, \sigma_{-1}\left(p_{k}\right)$ are all distinct. (In fact, $\sigma_{-1}\left(p_{1}\right)>$ $\left.\sigma_{-1}\left(p_{2}\right)>\cdots>\sigma_{-1}\left(p_{k}\right).\right)$

We can now prove the first half of Theorem 8.5, that the set of $n \in \mathbf{N}$ with $\sigma_{-1}(n) \leq u$ always possesses an asymptotic density:

Proposition 8.15. For each real $u$, the quantity $N(x, u) / x$ tends to a limit, say $D(u)$, as $x \rightarrow \infty$.

Proof. If $B_{1}<B_{2}$, then $\sigma_{-1}^{B_{1}}(n) \leq \sigma_{-1}^{B_{2}}(n)$ for each $n$, and so $D^{B_{1}}(u) \geq$ $D^{B_{2}}(u)$. Hence (for each fixed $\left.u\right) D^{B}(u)$ converges as $B \rightarrow \infty$ to $D^{*}(u):=$ $\inf _{B>0} D^{B}(u)$. We will prove that $N(x, u) / x \rightarrow D^{*}(u)$ as $x \rightarrow \infty$.

If $B>0$, then $\sigma_{-1}^{B}(n) \leq \sigma_{-1}(n)$ for every natural number $n$. Consequently, $\mathcal{N}(x, u) \subset \mathcal{N}^{B}(x, u)$ for all $x$. Thus

$$
\limsup _{x \rightarrow \infty} \frac{N(x, u)}{x} \leq \inf _{B>0}\left(\limsup _{x \rightarrow \infty} \frac{N^{B}(x, u)}{x}\right)=\inf _{B>0} D^{B}(u)=D^{*}(u) .
$$

We would like to establish the corresponding lower bound for the liminf of $N(x, u) / x$.

Let $\epsilon>0$. For a parameter $\delta>0$ to be specified shortly, put

$$
\mathcal{M}_{1}^{B}(x, u):=\left\{n \leq x: \sigma_{-1}^{B}(n) \leq u \text { and } u<\sigma_{-1}(n)<u+\delta\right\}
$$

and

$$
\mathcal{M}_{2}^{B}(x, u):=\left\{n \leq x: \sigma_{-1}^{B}(n) \leq u \text { and } \sigma_{-1}(n) \geq u+\delta\right\}
$$

and set $M_{i}^{B}(x, u):=\# \mathcal{M}_{i}^{B}(x, u)$. Then

$$
\begin{equation*}
\frac{N(x, u)}{x}=\frac{N^{B}(x, u)}{x}-\frac{M_{1}^{B}(x, u)}{x}-\frac{M_{2}^{B}(x, u)}{x} . \tag{8.13}
\end{equation*}
$$

If $\delta>0$ is small enough in terms of $\epsilon$, then $\lim \sup M_{1}^{B}(x, u) / x<\epsilon$ by Lemma 8.14. Having fixed such a $\delta$, notice that

$$
\begin{aligned}
M_{2}^{B}(x, u) & =\sum_{n \in \mathcal{M}_{2}^{B}(x, u)} 1 \leq \delta^{-1} \sum_{n \leq x}\left(\sigma_{-1}(n)-\sigma_{-1}^{B}(n)\right) \\
& =\delta^{-1} \sum_{\substack{d \leq x \\
p \mid d \text { for some } p>B}} \frac{1}{d} \sum_{\substack{n \leq x \\
d \backslash n}} 1 \leq \delta^{-1} x \sum_{d>B} d^{-2} \ll \delta^{-1} x / B .
\end{aligned}
$$

In particular, $\lim \sup M_{2}^{B}(x, u) / x$ tends to zero as $B \rightarrow \infty$. Letting first $x$ tend to infinity in (8.13) and then also $B$, we find

$$
\liminf \frac{N(x, u)}{x} \geq \lim _{B \rightarrow \infty} D^{B}(u)-\epsilon=D^{*}(u)-\epsilon
$$

Since $\epsilon>0$ is arbitrary, Proposition 8.15 follows.
The next proposition completes the proof of Theorem 8.3.
Proposition 8.16. If $D(u)$ is defined as in the statement of Proposition 8.15, then $D(u)$ defines a continuous function of $u$ on all of $\mathbf{R}$. Moreover, $D(1)=0$ and $D(u) \rightarrow 1$ as $u \rightarrow \infty$.

Proof. Clearly $D(u)$ is nondecreasing as a function of $u$. So if $u$ is an arbitrary real number, then for every real $\delta$ we have

$$
|D(u+\delta / 2)-D(u)| \leq|D(u+|\delta| / 2)-D(u-|\delta| / 2)|
$$

The right-hand side of this inequality represents the density of the set of $n$ for which $u-|\delta| / 2<\sigma_{-1}(n) \leq u+|\delta| / 2$, and this tends to zero with $\delta$ by Lemma 8.14. Hence $D$ is continuous at $u$.

It is clear that $D(1)=0$, since $\sigma_{-1}(n)>1$ except when $n=1$. Moreover, for all $x>0$,

$$
\sum_{n \leq x} \sigma_{-1}(n)=\sum_{d \leq x} \frac{1}{d} \sum_{\substack{n \leq x \\ d \mid n}} 1 \leq x \sum_{d} d^{-2}<2 x
$$

Thus, for each $u>0$, the number of $n \leq x$ with $\sigma_{-1}(n)>u$ is $<2 x / u$. Hence $1-D(u) \leq 2 / u$, so that $D(u) \geq 1-2 / u$. Since $D(u) \leq 1$ for all $u$, it follows that $D(u) \rightarrow 1$ as $u \rightarrow \infty$, as desired.

In Exercises 34 and 35 we outline a proof that $D(u)$ is strictly increasing for $u \geq 1$. In fact, that argument shows that $D(u)$ has an infinite right-sided derivative at every rational number $u$ of the form $\sigma_{-1}(n)$ (where $n \in \mathbf{N}$ ) while the set of such $u$ is dense in $[1, \infty)$. Erdős has proved $[\mathbf{E r d} 39]$ the curious result that $D^{\prime}(u)=0$ for all $u$ outside of a set of measure zero.

## 5. Sociable numbers

5.1. A theorem on the local behavior of aliquot sequences. One way to disprove the Catalan-Dickson conjecture mentioned in this chapter's introduction would be to produce a natural number $n$ for which the sequence $\left\{s_{j}(n)\right\}_{j=0}^{\infty}$ is strictly increasing. It seems unlikely that such an $n$ exists. However, in 1975 Lenstra [Len75] showed that for each fixed $K$, there are infinitely many natural numbers $n$ with

$$
\begin{equation*}
n<s(n)<s_{2}(n)<\cdots<s_{K+1}(n) . \tag{8.14}
\end{equation*}
$$

Actually (8.14) is more common than one might expect: In 1976, Erdős showed [Erd76] that for each fixed $K$, (8.14) holds for almost all abundant numbers $n$. In other words, if $n$ increases once when $s$ is applied, then almost surely $n$ increases $K+1$ times. Erdős deduced this result from the following theorem, which is of independent interest:

Theorem 8.17. Let $K$ be a natural number, and let $\epsilon>0$. For almost all natural numbers $n$,

$$
\frac{s_{k+1}(n)}{s_{k}(n)}>\frac{s(n)}{n}-\epsilon
$$

for all $1 \leq k \leq K$.
Before proceeding to the proof of Theorem 8.17, let us see how to derive the stated consequence:

Corollary 8.18. For each fixed $k$, the set of abundant numbers $n$ for which (8.14) fails has asymptotic density zero.

Proof. Let $\epsilon>0$. Using the continuity of the distribution function $D(u)$ of Theorem 8.5, choose a small $\delta>0$ with $D(2+\delta)<D(2)+\epsilon$. Suppose $n$ is abundant but that (8.14) fails. If $\sigma(n) / n \leq 2+\delta$, then $n$ belongs to a set of density $D(2+\delta)-D(2)<\epsilon$. Now suppose that $\sigma(n) / n>2+\delta$. By Theorem 8.17, unless $n$ belongs to a certain set of density zero,

$$
s_{k+1}(n) / s_{k}(n)>s(n) / n-\delta / 2>(1+\delta)-\delta / 2>1
$$

for all $1 \leq k \leq K$, and so (8.14) holds.
So the set of abundant counterexamples to (8.14) has upper density less than $\epsilon$. Since $\epsilon>0$ was arbitrary, the corollary follows.

The proof of Theorem 8.17 requires a preliminary technical lemma.
Lemma 8.19. Let $K$ and $M$ be integers with $K \geq 0$ and $M \geq 1$. Then the following is true for almost all natural numbers $n$ : There are primes $p_{0}, p_{1}, \ldots, p_{K}$ for which

$$
\begin{equation*}
p_{i} \| n \quad \text { for each } i=0,1,2, \ldots, K, \tag{8.15}
\end{equation*}
$$

and

$$
\begin{equation*}
p_{0} \equiv-1(\bmod M), \quad \text { and } \quad p_{i+1} \equiv-1\left(\bmod p_{i}^{2}\right) \quad \text { for all } \quad 0 \leq i<K . \tag{8.16}
\end{equation*}
$$

Proof. The lemma is a consequence of the following assertion, which we prove by induction on $K$ : For each nonnegative integer $K$, each $M \in \mathbf{N}$, and each $\epsilon>0$, there is a number $B$ with the property that for all $n$ outside of a set of upper density $<\epsilon$, one can find primes $p_{0}, \ldots, p_{K} \leq B$ satisfying both (8.15) and (8.16). When $K=0$, this statement follows immediately from Lemma 8.13, applied with

$$
\mathcal{P}:=\{p \equiv-1(\bmod M)\} .
$$

(Note that $\sum_{p \in \mathcal{P}} p^{-1}$ diverges by the results of Chapter 4.)
Now suppose the statement is known to hold for a certain integer $K \geq 0$. If $M \in \mathbf{N}$ and $\epsilon>0$ are given, the induction hypothesis permits us to choose a number $B_{0}$ with the property that for all $n$ outside of a set $E_{0}$ (say) of upper density $<\epsilon / 2$, there are primes $p_{0}, \ldots, p_{K} \leq B_{0}$ satisfying (8.15) and (8.16). Let $R:=\left(\prod_{p \leq B_{0}} p\right)^{2}$ and apply Lemma 8.13 with $\mathcal{P}:=\{p \equiv$ $-1(\bmod R)\}$. We find that for a suitable choice of $z$, all $n$ outside of a set $E_{1}$ (say) of upper density $<\epsilon / 2$ have an exact prime divisor $p_{K+1} \equiv$ $-1(\bmod R)$ with $p_{K+1} \leq z$. But then if $n$ lies outside $E_{0} \cup E_{1}$, the primes $p_{1}, \ldots, p_{K+1}$ satisfy (8.15) and (8.16) with $K$ replaced by $K+1$. Since $E_{0} \cup E_{1}$ has upper density $<\epsilon$, we obtain the $(K+1)$-case of the assertion with $B=\max \left\{B_{0}, z\right\}$.

Proof of Theorem 8.17. Let $B$ be an arbitrary natural number, and put $M:=\left(\prod_{p \leq B} p\right)^{B}$. We claim that for almost all $n$, the number $M$ divides $\sigma\left(s_{i}(n)\right)$ for each $0 \leq i \leq K$.

The proof of the claim starts with the observation that by Lemma 8.19, for almost all $n$ there are primes $p_{0}, \ldots, p_{K}$ satisfying (8.15) and (8.16). Then for each $0 \leq i<K$, we have

$$
p_{i}^{2}\left|\sigma\left(p_{i+1}\right)\right| \sigma(n), \quad \text { so that since } p_{i} \| n, \text { we have } \quad p_{i} \| \sigma(n)-n=s(n)
$$

Thus $p_{0}, \ldots, p_{K-1}$ exactly divide $s(n)$. We can repeat the argument with $n$ replaced by $s(n)$ to see that $s_{2}(n)$ is exactly divisible by $p_{0}, \ldots, p_{K-2}$. Continuing in the same manner, we find that $s_{i}(n)$ is exactly divisible by
$p_{0}, \ldots, p_{K-i}$, for each $0 \leq i \leq K$. In particular, $p_{0}$ exactly divides each of $n, s(n), \ldots, s_{K}(n)$. Thus

$$
M\left|\sigma\left(p_{0}\right)\right| \sigma\left(s_{i}(n)\right) \quad \text { for all } 0 \leq i \leq K
$$

as we originally claimed.
So at the cost of throwing away a set of density zero, we may assume that the claim holds for $n$. As a consequence, for each $0<k \leq K+1$, we have

$$
\begin{equation*}
s_{k}(n)=\sigma\left(s_{k-1}(n)\right)-s_{k-1}(n) \equiv-s_{k-1}(n)(\bmod M) . \tag{8.17}
\end{equation*}
$$

For each $0 \leq i \leq K$, write $s_{i}(n)=m_{i} n_{i}$, where $\operatorname{gcd}\left(m_{i}, n_{i}\right)=1$ and every prime divisor of $n_{i}$ is at most $B$. (So $n_{i}$ is the $B$-smooth part of $s_{i}(n)$.) We claim that for all $n$ outside of a set of upper density $o(1)$, we have

$$
\begin{equation*}
n_{0}=n_{1}=\cdots=n_{K} ; \tag{8.18}
\end{equation*}
$$

here and below, $o(1)$ denotes a quantity that tends to zero as $B \rightarrow \infty$. For the proof, suppose (8.18) fails, so that $n_{i} \neq n_{i+1}$ for some $0 \leq i<K$. Writing

$$
s_{i+1}(n)=\sigma\left(s_{i}(n)\right)-s_{i}(n),
$$

we see that $n_{i} \neq n_{i+1}$ implies that there is a prime $p \leq B$ which divides $s_{i}(n)$ to at least as high a power as it divides $\sigma\left(s_{i}(n)\right)$. Since $\sigma\left(s_{i}(n)\right)$ is divisible by $M$, and hence by $p^{B}$, it must be that $p^{B}$ divides $s_{i}(n)$. But then by repeated application of (8.17) (starting with $k=i$ ), we find that $p^{B}$ divides $s_{0}(n)=n$. But the upper density of the set of $n$ divisible by $p^{B}$ for some $p \leq B$ is bounded by $\sum_{p \leq B} p^{-B}$, which is $o(1)$.

So, excepting a set of upper density $o(1)$, we may suppose that (8.18) holds. Then for each $1 \leq k \leq K$,

$$
\begin{aligned}
\frac{s(n)}{n}-\frac{s_{k+1}(n)}{s_{k}(n)} & =\frac{\sigma(n)}{n}-\frac{\sigma\left(s_{k}(n)\right)}{s_{k}(n)} \\
& =\frac{\sigma\left(n_{0}\right)}{n_{0}}\left(\frac{\sigma\left(m_{0}\right)}{m_{0}}-\frac{\sigma\left(m_{k}\right)}{m_{k}}\right) \\
& \leq \frac{\sigma\left(n_{0}\right)}{n_{0}}\left(\frac{\sigma\left(m_{0}\right)}{m_{0}}-1\right) .
\end{aligned}
$$

Now $\sigma\left(n_{0}\right) / n_{0}=\sigma_{-1}\left(n_{0}\right) \leq \sigma_{-1}(n) ;$ moreover, $\sigma_{-1}(n) \leq B^{1 / 2}$ for all $n$ outside of a set of density $o(1)$, by the latter half of Lemma 8.16. We claim that we also have

$$
\frac{\sigma\left(m_{0}\right)}{m_{0}}-1 \leq \frac{1}{B^{3 / 4}}
$$

for all $n$ outside of a set of upper density $o(1)$. Once this claim is established, we will have shown that for all $n$ outside of a set of upper density $o(1)$,

$$
\frac{s(n)}{n}-\frac{s_{k+1}(n)}{s_{k}(n)} \leq \frac{B^{1 / 2}}{B^{3 / 4}}=\frac{1}{B^{1 / 4}}=o(1) \quad \text { for all } 1 \leq k \leq K
$$

and Theorem 8.17 follows upon letting $B \rightarrow \infty$.
To prove this last claim, notice that

$$
\frac{\sigma\left(m_{0}\right)}{m_{0}}-1=\sum_{\substack{d|n \\ p| d \Rightarrow p>B \\ d>1}} \frac{1}{d},
$$

so that the number of $n \leq x$ with $\sigma\left(m_{0}\right) / m_{0}-1>B^{-3 / 4}$ is at most

$$
\begin{aligned}
B^{3 / 4} \sum_{n \leq x} \sum_{\substack{1<d|n \\
p| d \Rightarrow p>B}} \frac{1}{d} & =B^{3 / 4} \sum_{\substack{1<d \leq x \\
p \mid d \Rightarrow p>B}} \frac{1}{d} \sum_{\substack{n \leq x \\
d \mid n}} 1 \\
& \leq B^{3 / 4} \sum_{d>B} \frac{x}{d^{2}} \leq B^{-1 / 4} x .
\end{aligned}
$$

Thus the set of such $n$ has upper density $\leq B^{-1 / 4}=o(1)$, as desired.
5.2. An application to sociable numbers. Theorem 8.6, which asserts that the set of sociable numbers of order $k$ has density zero for each fixed $k$, is almost immediate from Corollary 8.18. Indeed, fix a natural number $k>1$. (When $k=1$, we have already seen that the sociable numbers of order $k$ - i.e., the perfect numbers - comprise a set of density zero.) Let $A(x)$ be the number of sociable $n \leq x$ of order $k$, and let $A^{\prime}(x)$ be the number of $n \leq x$ which are the smallest member of some sociable $k$-cycle. Then $A(x) \leq k A^{\prime}(x)$. So to show that $A(x)=o(x)$, it is enough to show that $A^{\prime}(x)=o(x)$. But this is clear from Corollary 8.18, since if $n$ is the smallest member of a sociable $k$-cycle, then $n<s(n)$ (i.e., $n$ is abundant), but we do not have

$$
n<s(n)<s_{2}(n)<\ldots<s_{k}(n),
$$

since $n=s_{k}(n)$.
Remark. By making the argument above explicit when $k=2$, Erdős \& Rieger ([Rie73], [ER75]) showed that the number of amicable $n \leq x$ is $\ll x / \log \log \log x$. Let $\log _{1} x:=\max \{1, \log x\}$ and for $k>1$, define $\log _{k} x:=$ $\max \left\{1, \log \left(\log _{k-1} x\right)\right\}$. For general $k$, the Erdős-Rieger method shows that there are $<_{k} x / \log _{r} x$ sociable numbers of order $k$ not exceeding $x$, where $r$ grows linearly with $k$ (e.g., $r=3 k$ is permissible). In [KPP09], it is proved
that the number of sociable $n \leq x$ of order $k$ is at most

$$
k\left(2 \log _{4} x\right)^{k} \frac{x}{\exp \left((1+o(1)) \sqrt{\log _{3} x \log _{4} x}\right)}
$$

where the $o(1)$ term tends to zero as $x \rightarrow \infty$, and the estimate is uniform in $k \geq 1$. Moreover, for odd $k$, one can do a bit better; in this case the count is

$$
\ll k\left(2 \log _{4} x\right)^{k} \frac{x}{\sqrt{\log _{2} x \log _{3} x}},
$$

where the implied constant is absolute.

## Notes

The first chapter of Dickson's History of the Theory of Numbers [Dic66] is a thorough compendium of results on perfect numbers and related matters, covering antiquity to the early twentieth century. Many of the more recent results (up to about 2003) are catalogued in the two-volume Handbook of Number Theory; see, in particular, $[\mathbf{S C 0 4}$, Chapter 3] and $[\mathbf{S M C 0 6}$, Chapter 1].

Theorems 8.3 and 8.4 can both be generalized. In fact, what Wirsing actually shows in [Wir59] is that for any $\alpha$, the number of $n \leq x$ with

$$
\begin{equation*}
\sigma(n) / n=\alpha \tag{8.19}
\end{equation*}
$$

is at most $x^{W / \log \log x}$, for an absolute constant $W>0$ (and all $x \geq 3$ ). The complete independence from $\alpha$ of this upper bound is frequently useful in applications. As for Dickson's finiteness theorem, the following elegant generalization was proved by Kanold [Kan56]: Call a solution $n$ to (8.19) primitive if $n$ does not have a unitary divisor which is an even perfect number. ${ }^{3}$ For each $\alpha \in \mathbf{Q}$ and $k \in \mathbf{N}$, there are only finitely many primitive solutions $n$ to (8.19) with exactly $k$ distinct prime factors. In [Pom77a], Pomerance shows how Baker's estimates for linear forms in logarithms can be used to obtain an effective version of Kanold's result. Borho [Bor74a, Bor74b] and Artjuhov $[\mathbf{A r t 7 5}]$ have obtained results for amicable pairs which are cognate to Dickson's theorem.

In the theory of probability, a function $F: \mathbf{R} \rightarrow \mathbf{R}$ is called a distribution function if $F$ is nondecreasing, right-continuous,

$$
\lim _{u \rightarrow-\infty} F(u)=0 \quad \text { and } \quad \lim _{u \rightarrow \infty} F(u)=1
$$

We say that an arithmetic function $f$ has a distribution function if there is a distribution function $D_{f}$ (say) with the property that

$$
\lim _{x \rightarrow \infty} \frac{\#\{n \leq x: f(n) \leq u\}}{x}=D_{f}(u)
$$

[^6]whenever $u$ is a point of continuity of $D_{f}$. The result of Davenport recorded in Theorem 8.5 is an early precursor of the following theorem of Erdős ([Erd35a, Erd37, Erd38]) \& Wintner [EW39]:
$\star$ Theorem 8.20. A real-valued additive arithmetic function $f(n)$ has a distribution function if and only if all of the three series
$$
\sum_{|f(p)|>1} \frac{1}{p}, \quad \sum_{|f(p)| \leq 1} \frac{f(p)}{p}, \quad \sum_{|f(p)| \leq 1} \frac{f(p)^{2}}{p}
$$
converge. If all three series converge, then the distribution function of $f$ is continuous if and only $\sum_{f(p) \neq 0} p^{-1}$ diverges.

Of course $\sigma(n) / n$ is multiplicative, not additive, but one can recover Theorem 8.5 by applying Theorem 8.20 to $\log (\sigma(n) / n)$. The Erdős-Wintner result can be considered the first general theorem in the subject that has come to be known as "probabilistic number theory".

Theorem 8.17 says that for most natural numbers $n$, the aliquot sequence $n, s(n), s(s(n)), \ldots$ initially grows almost as fast as a geometric progression with common ratio $s(n) / n$. While technical, our proof from $\S 5$ can be summarized neatly in one sentence: For most $n$, the first few terms of the aliquot sequence at $n$ have all of the same small prime factors, while for most $m$, the ratio $\sigma(m) / m$ is "nearly determined" by the small prime factors of $m$. This summary might lead one to expect that one should also have the statement analogous to Theorem 8.17 where the inequality points in the opposite direction. This was conjectured by Erdős [Erd76]:

Conjecture 8.21. Let $K$ be a natural number, and let $\epsilon>0$. For almost all natural numbers $n$,

$$
\frac{s_{k+1}(n)}{s_{k}(n)}<\frac{s(n)}{n}+\epsilon
$$

for all $1 \leq k \leq K$.
This conjecture has proved surprisingly difficult and remains open in general. For a proof when $K=1$, see the paper [EGPS90] of Erdős et al.

## Exercises

1. (Lucas) The digital root of a natural number $n$ is defined by summing the (decimal) digits of $n$, then the digits of the result, then the digits of the new result, etc., until reaching a single digit. Prove that every even perfect number $n>6$ has digital root 1 .
2. Identify the flaw in the following "proof" that all perfect numbers $n$ are even: Starting with $2 n=\sum_{d \mid n} d$, we can apply Möbius inversion to find that

$$
n=\sum_{d \mid n} \mu(n / d)(2 d)=2\left(\sum_{d \mid n} \mu(n / d) d\right)
$$

which is visibly even.
3. (Ewell $[\mathbf{E w e 8 0}]$ ) Suppose that $n$ is an odd perfect number. Write $n=$ $p^{e} \prod_{i=1}^{r} p_{i}^{2 e_{i}} \prod_{j=1}^{s} q_{j}^{2 f_{j}}$, where $p, p_{1}, \ldots, p_{r}, q_{1}, \ldots, q_{s}$ are distinct primes, $p \equiv e \equiv 1(\bmod 4)$, each $p_{i} \equiv 1(\bmod 4)$, and each $q_{j} \equiv 3(\bmod 4)$. Show that $p \equiv e(\bmod 8)$ precisely when there are an even number of odd $e_{i}$.
4. (Starni [Sta91]) Suppose that $n$ is an odd perfect number. Write $n=$ $p^{e} \prod_{i=1}^{k} p_{i}^{2 e_{i}}$, as in Euler's theorem (Theorem 8.2).
(a) Show that if $p_{i} \equiv 3(\bmod 4)$ for all $1 \leq i \leq k$, then $\frac{1}{2} \sigma\left(p^{e}\right)$ is composite.
(b) Show that if $p_{i} \equiv 1(\bmod 4)$ for all $1 \leq i \leq k$, then $p \equiv e(\bmod 8)$.
5. (Starni [Sta93]) Let $n$ be an odd perfect number, say $n=p^{e} m^{2}$, where $\operatorname{gcd}(p, m)=1$ and $p \equiv e \equiv 1(\bmod 4)$. Show that if $e+2$ is a prime which does not divide $p-1$, then $e+2$ divides $m^{2}$. For example, if $13^{17} m^{2}$ is perfect (with $13 \nmid m$ ), then 19 divides $m$.
6. (Slowak [Slo99]) Let $n$ be an odd perfect number, say $n=p^{e} m^{2}$, where $\operatorname{gcd}(p, m)=1$ and $p \equiv e \equiv 1(\bmod 4)$. Show that $p^{e}$ is a proper divisor of $\sigma\left(m^{2}\right)$.
7. (Touchard [Tou53]) Show that if $n$ is an odd perfect number, then either $n \equiv 1(\bmod 12)$ or $n \equiv 9(\bmod 36)$.
8. (Luca [Luc99]) Prove that two consecutive natural numbers cannot both be perfect.
 the Euler-Mascheroni constant.

Remark. A handsome theorem of Robin [Rob84] asserts that the Riemann Hypothesis (see p. 105) holds if and only if $\sigma(n)<e^{\gamma} n \log \log n$ for all $n>5040$.
10. (Salié [Sal53]) Let $n$ be an abundant or perfect number with $k$ distinct prime factors, and let $q$ be its least prime divisor. Let $q^{\prime}$ be the $k$ th prime exceeding $q$. Observing that

$$
2 \leq \frac{\sigma(n)}{n}<\prod_{p \mid n} \frac{p}{p-1} \leq \prod_{q \leq p<q^{\prime}}\left(1-\frac{1}{p}\right)^{-1},
$$

deduce from Mertens' theorem that $q \ll \sqrt{k \log k}$, where the implied constant is absolute. Some related results can be found in the paper [Nor61] of Norton.
11. (Yamada [Yam]) Let $E$ be a finite, nonempty set of natural numbers. Let $n$ be an odd perfect number, and suppose that every even exponent appearing in the canonical prime factorization of $n$ belongs to the set $\{2 e: e \in E\}$. Put $\mathcal{Q}:=\left\{q\right.$ prime $\left.: q \equiv 1\left(\bmod \prod_{e \in E}(2 e+1)\right), q \nmid 2 n\right\}$.
(a) Suppose $e \in E$ and that $p^{2 e} \| n$ for the prime $p$. Show that for each $q \in \mathcal{Q}$,

$$
1+p+p^{2}+\cdots+p^{2 e} \not \equiv \equiv 0(\bmod q) .
$$

(b) Show that the polynomial $1+T+T^{2}+\cdots+T^{2 e}$ has exactly $2 e$ distinct roots modulo $q$, for each $q \in \mathcal{Q}$.
(c) Let $x \geq 3$. Show that for each $e \in E$, the number of primes $p \leq x$ for which $p^{2 e} \| n$ is $\ll x /(\log x)^{1+c}$, where $c>0$, and where both $c$ and the implied constant depend only on $E$ (and not on $n$ ). Hint: Apply the Brun-Hooley sieve.
(d) Let $n^{\prime}$ be the product of the prime powers with even exponent which exactly divide $n$. Show that $\sigma\left(n^{\prime}\right) / n^{\prime} \geq 8 / 5$. Deduce that for some choice of $e \in E$,

$$
(8 / 5)^{1 / \# E} \leq \prod_{p: p^{2 e} \| n}\left(1+\frac{1}{p-1}\right) .
$$

(e) Combining the results of (c) and (d), show that if $p$ is the smallest prime appearing to an even exponent in $n$, then $p$ is bounded above by a constant depending only on $E$.
12. (Anderson [And74]) Show that if $\sigma(n) / n=5 / 3$, then $n$ is coprime to 10. Deduce that in this case $5 n$ is an odd perfect number.
13. (Anderson, ibid.)
(a) Suppose $1 \leq b \leq a<\sigma(b)$ and $\operatorname{gcd}(a, b)=1$. Prove that the rational number $a / b$ is not of the form $\sigma(n) / n$ for any $n \in \mathbf{N}$.
(b) Show that the rational numbers not of the form $\sigma(n) / n$ are dense in $[1, \infty)$.
Further results related to those of Exercises 12 and 13 may be found in [Wei00], [Hol06], and [SH07].
14. Call the natural number $n$ superperfect if $\sigma(\sigma(n))=2 n$.
(a) (Suryanarayana [Sur69]) Show that if $n$ is an even superperfect number, then $n=2^{k}$ for some $k \in \mathbf{N}$ for which $2^{k+1}-1$ is prime. Conversely, show that any $n$ of this form is superperfect.
(b) (Kanold [Kan69b]) Show that if $n$ is an odd superperfect number, then $n$ is a perfect square. (No examples of odd superperfect numbers are known.)
15. (Makowski [Mąk62]) Show that 28 is the only even perfect number of the form $m^{3}+1$ and the only even perfect number of the form $m^{m}+1$.
16. (Wall [Wal81])
(a) Prove that for every $k \in \mathbf{N}$, there are infinitely many blocks of $k$ consecutive natural numbers all of which are abundant.
(b) Show that there are infinitely many blocks of 5 consecutive numbers all of which are deficient and that 5 cannot be replaced with any larger number.
17. Show that every sufficiently large natural number can be written as a sum of two abundant numbers and as a sum of two deficient numbers.
18. (Pomerance [unpublished], te Riele [tR76, Chapter 7]) The Dedekind $\psi$ function is defined by setting $\psi(n):=n \prod_{p \mid n}(1+1 / p)$ for every natural number $n$. (Thus $\psi(n) \leq \sigma(n)$ for all $n$, with equality precisely when $n$ is squarefree.) Show that the analogue of the Catalan-Dickson conjecture fails for $s^{*}(n):=\psi(n)-n$. That is, there are natural numbers $n$ for which the sequence $n, s^{*}(n), s^{*}\left(s^{*}(n)\right), \ldots$ is unbounded. Hint: Try $n=318$.
19. (Alaoglu \& Erdős [AE44]) Prove that for each fixed $\epsilon>0$, the inequality $\varphi(\sigma(n))<\epsilon n$ holds on a set of $n$ of density 1 .
20. (Kanold [Kan69a], see also Borho [Bor70]) It is not known whether an amicable number can possess only a single prime factor (and so be a prime or prime power). Show that the number of amicable numbers of this type not exceeding $x$ is $O_{\epsilon}\left((\log x)^{1+\epsilon}\right)$ as $x \rightarrow \infty$, for each $\epsilon>0$.
21. (a) (Dirichlet) Show that $\frac{\sigma(n)}{n}$ has mean value $\frac{\pi^{2}}{6}$. In other words, prove that $\frac{1}{x} \sum_{n \leq x} \frac{\sigma(n)}{n} \rightarrow \frac{\pi^{2}}{6}$ as $x \rightarrow \infty$.
(b) (Erdős $[\mathbf{E r d} 51]$ ) Prove that $\frac{\sigma\left(2^{n}-1\right)}{2^{n}-1}$ possesses a (finite) mean value. Hint: Use the result of Exercise 6.34(c).
22. (Bojanić $[\mathbf{B o j} 54])$ Show that $\frac{\sigma\left(2^{p}-1\right)}{2^{p}-1} \rightarrow 1$ as $p \rightarrow \infty$ through prime values.
23. (Luca [Luc00a]) Let $F_{m}=2^{2^{m}}+1$ be the $m$ th Fermat number. Show that $s\left(F_{m}\right) \ll m F_{m} / 2^{m}$ for $m \geq 1$. Combining this with the result of Exercise 9, prove that only finitely many Fermat numbers $F_{m}$ are members of an amicable pair. (With a bit of extra work, one can show that there are no such numbers.)
24. (Luca [Luc06, Problem 171]) Call the natural number $n$ multiply perfect if $n$ divides $\sigma(n)$. Show that for each fixed $B$, there are only finitely many multiply perfect numbers all of whose prime factors are bounded by $B$.
25. (Pomerance [Pom93]) Prove that $n!$ is multiply perfect for only finitely many $n$. (It can be shown that $n=1,3$, and 5 are the only such $n$.) Hint: One argument starts by showing that $v \ll n / \log n$ as $n \rightarrow \infty$, where $v=v(n)$ is defined by the relation $2^{v} \| \sigma(n!)$.

Remark. A plausible strengthening of the result of this exercise was suggested by Erdős: It is not hard to check that as $n \rightarrow \infty$, we have $\sigma(n!) / n!\sim e^{\gamma} \log \log n!($ cf. Exercise 9). Erdős's conjecture is that for each $\epsilon>0$, there are only finitely many multiply perfect $m$ with $\sigma(m) / m>\epsilon \log \log m$.
26. A natural number $m$ is called untouchable if it is not of the form $s(n)$ for any $n \in \mathbf{N}$. The sequence of untouchable numbers begins $2,5,52,88,96,120,124,146, \ldots$.
(a) Prove that $s(n)>\sqrt{n}$ for every composite number $n$. Using this inequality (or not) check that 2 and 5 are both untouchable.
(b) Show that if every even number $m \geq 8$ is the sum of two distinct primes (a conjecture strengthening Goldbach's), then 5 is the only odd untouchable number.
27. (Continuation; Erdős [Erd73], see also [BL05], [tR76]) We now show that a positive proportion of natural numbers are untouchable.

Let $M$ be a fixed even natural number. We consider the inequality

$$
\begin{equation*}
s(n) \leq x, \quad \text { with the constraint } \quad s(n) \equiv 0(\bmod M) . \tag{8.20}
\end{equation*}
$$

(a) Show that the number of odd $n$ for which (8.20) holds is $\ll x / \log x$ as $x \rightarrow \infty$. Hint: $\sigma(n)$ is odd only if $n$ is a perfect square.
(b) Show that the number of solutions to (8.20) in even numbers $n$ not divisible by $M$ is $o(x)$.
(c) Show that the number of solutions to (8.20) in numbers $n$ which are divisible by $M$ is at most $\alpha x$, where $\alpha:=(\sigma(M)-M)^{-1}$.
(d) Combining the results of (a)-(c), deduce that the number of solutions to $(8.20)$ is $\leq(\alpha+o(1)) x$.
(e) Taking $M=12$, show that at least $\left(\frac{1}{48}+o(1)\right) x$ natural numbers $n \leq x$ are untouchable.
28. Show that for each fixed natural number $k$ and rational number $\alpha$, the set of natural numbers $n$ with $s_{k}(n)=\alpha n$ has density zero.
29. (Banks et al. [BFPS04]) Show that there are infinitely many natural numbers $n$ for which $\sigma(n)$ is a perfect square. This had been conjectured by Sierpiński [Sie88, p. 179]. Hint: View the group $\mathbf{Q}^{\times} /\left(\mathbf{Q}^{\times}\right)^{2}$ as an $\mathbf{F}_{2}$-vector space, with a basis given by the images of -1 and the rational
primes. Show that there are many linear dependencies in $\mathbf{Q}^{\times} /\left(\mathbf{Q}^{\times}\right)^{2}$ among the shifted primes $p+1$. Exercise 6.25 will prove useful.
30. (Pomerance $[\mathbf{P o m 7 7 b}])$ Show that $\tau(n)$ divides $\sigma(n)$ for almost all natural numbers $n$. That is, the arithmetic mean of the divisors of $n$ is almost always an integer.

Remark. In [BEPS81], it is shown that the number of exceptional $n \leq x$ is $x \exp (-(2+o(1)) \sqrt{\log 2} \sqrt{\log \log x})$ as $x \rightarrow \infty$.
31. (Adapted from [Luc06, Problem 148]) Fill in the details in the following proof that the arithmetic mean of the distinct prime divisors of $n$ is almost never an integer (i.e., is an integer only on a set of density zero):

Let $n \leq x$. We can assume that the largest prime divisor $P(n)$ of $n$ satisfies $P(n)>y$, where $y:=x^{1 / \log \log \log x}$, since the $n \leq x$ for which this fails make up a set of size $o(x)$ by Exercise 6.19. Write $n=P m$, where $P=P(n)$. We can further assume $P \nmid m$, since otherwise $n$ has a large square divisor, and such $n$ are also rare. Finally, Exercise 3.23 allows us to assume that $\omega(n) \in\left[\log \log x-(\log \log x)^{2 / 3}, \log \log x+(\log \log x)^{2 / 3}\right]$. If the average of the prime divisors of $n$ is an integer, this forces $P$ to lie in a residue class, modulo $\omega(m)+1$, determined entirely by $m$. We now consider the number of possibilities for $P$ corresponding to a given $m \leq x / y$. With $k:=\omega(m)+1$, the number of suitable $P \leq x / m$ is, by the Brun-Titchmarsh inequality (Exercise 6.21),

$$
\begin{aligned}
\ll \frac{x / m}{\varphi(k) \log (x /(m k))} & \ll \frac{x \log \log \log x}{m \varphi(k) \log x} \\
& \ll \frac{x(\log \log \log x)(\log \log \log \log x)}{m(\log \log x)(\log x)},
\end{aligned}
$$

where we use that $x / m \geq y$, that $k \approx \log \log x$, and that $\varphi(r) \gg r / \log \log r$ for all $r \geq 3$ (cf. Exercise 9). The result follows upon summing over the possibilities for $m$.

Remark. For strengthenings of this result, see the papers of Banks et al. [BGLS05] and Kátai [Kát07].
32. (Erdős [Erd46]) In this exercise we investigate the rate at which $D(u) \rightarrow$ 1 as $u \rightarrow \infty$, where $D(u)$ is the function of Theorem 8.5. We show that the density $1-D(u)$ of those $n$ for which $\sigma(n) / n>u$ is

$$
\begin{equation*}
1 / \exp \left(\exp \left(\left(e^{-\gamma}+o(1)\right) u\right)\right) . \tag{8.21}
\end{equation*}
$$

The bulk of the proof (parts (a)-(d)) concerns the upper bound. Actually we prove a somewhat stronger result, namely that (8.21) is an upper bound for the upper density of the set of $n$ with

$$
\begin{equation*}
\prod_{p \mid n}(1-1 / p)^{-1}>u \tag{8.22}
\end{equation*}
$$

(Notice that the left-hand side of (8.22) majorizes $\sigma(n) / n$.)
(a) With $p_{i}$ denoting the $i$ th prime, let $k=k(u)$ be the smallest natural number with $\prod_{i=1}^{k}\left(1-1 / p_{i}\right)^{-1}>u$. Prove that $\log p_{k} \sim e^{-\gamma} u$ as $u \rightarrow \infty$.
(b) Divide the solutions $n$ of (8.22) into two classes:
(i) $n$ has at least $r:=\lfloor k / 2\rfloor$ prime factors not exceeding $4 p_{k}$,
(ii) all other solutions to (8.22).

Show that class (i) has upper density at most $2^{\pi\left(4 p_{k}\right)} / \prod_{p \leq p_{r}} p$. Use (a) and the prime number theorem to verify that this bound has the form (8.21). Thus we may focus attention on class (ii).
(c) Use the minimality of $k$ to show that if $n$ is a solution to (8.22) belonging to class (ii), then (for large $u$ )

$$
\prod_{\substack{p \mid n \\ p>4 p_{k}}}(1-1 / p) \leq \prod_{j=r+1}^{k-1}\left(1-1 / p_{j}\right)<1-\frac{1}{4 \log k}
$$

Deduce that for some natural number $j$,

$$
\sum_{\substack{p \mid n \\ 4^{j} p_{k}<p \leq 4^{j+1} p_{k}}} \frac{1}{p}>\frac{1}{2^{j}} \cdot \frac{1}{4 \log k},
$$

so that $n$ is divisible by at least $N_{j}:=\left\lceil 2^{j} \frac{p_{k}}{4 \log k}\right\rceil$ distinct primes from the interval $I_{j}:=\left(4^{j} p_{k}, 4^{j+1} p_{k}\right]$.
(d) Conclude that the upper density of class (ii) is, for large $u$, bounded above by $\sum_{j=1}^{\infty} \frac{1}{N_{j}!}\left(\sum_{p \in I_{j}} 1 / p\right)^{N_{j}}$. Check that this is, in turn, bounded above by an expression of the shape (8.21).
(e) It remains only to prove that (8.21) is a lower bound for the density of the set of $n$ with $\sigma(n) / n>u$. To accomplish this, construct a number

$$
n_{0} \leq \exp \left(\exp \left(\left(e^{-\gamma}+o(1)\right) u\right)\right)
$$

with $\sigma\left(n_{0}\right) / n_{0}>u$, and observe that $\sigma(n) / n>u$ whenever $n_{0}$ divides $n$. (Cf. Exercise 9.)
33. (Erdős, ibid.) Now we consider the decay of $D(u)$ to 0 as $u$ tends down to 1 . We show that the set of $n$ with $\sigma(n) / n \leq 1+\epsilon$ has density $\sim e^{-\gamma} / \log \left(\epsilon^{-1}\right)$ as $\epsilon \downarrow 0$.
(a) Let $A_{\epsilon}$ be the set of natural numbers with no prime factor $<\epsilon^{-1}$. Show that if $\sigma(n) / n \leq 1+\epsilon$, then $n \in A_{\epsilon}$.
(b) Show that $A_{\epsilon}$ has density $(1+o(1)) e^{-\gamma} / \log \left(\epsilon^{-1}\right)$ as $\epsilon \downarrow 0$.
(c) Prove that if $\epsilon$ is sufficiently small, then the following holds: If $n \in A_{\epsilon}$ but $\sigma(n) / n>1+\epsilon$, then for some natural number $j$, $n$ has at least $j$ distinct prime factors from the interval $I_{j}:=$ $\left[4^{j-1} \epsilon^{-1}, 4^{j} \epsilon^{-1}\right)$.
(d) For each natural number $j$, let $E_{j}$ be the set of $n \in A_{\epsilon}$ with at least $j$ distinct prime factors from $I_{j}$. Show that $E_{j}$ has upper density at most

$$
(1+o(1)) \frac{1}{j!} \frac{e^{-\gamma}}{\log (1 / \epsilon)}\left(\sum_{p \in I_{j}} \frac{1}{p}\right)^{j} .
$$

(e) Show that $\bigcup_{j \geq 1} E_{j}$ has upper density at most

$$
(1+o(1)) \frac{e^{-\gamma}}{\log (1 / \epsilon)} \sum_{j=1}^{\infty} \frac{1}{j!}\left(\sum_{p \in I_{j}} \frac{1}{p}\right)^{j} .
$$

(f) Complete the proof by showing that the sum in (e) tends to zero as $\epsilon \downarrow 0$.
34. Suppose that $f$ is a nonnegative-valued additive function for which
(i) $f(p) \rightarrow 0$ as $p \rightarrow \infty$,
(ii) $\sum_{p} f(p)$ diverges.

Show that the image of $f$ is dense in $[0, \infty)$. Taking $f(n):=\log \frac{\sigma(n)}{n}$, conclude that the set of rational numbers of the form $\sigma(n) / n$ is dense in $[1, \infty)$.
35. (a) Use the result of Exercise 33 to show that the function $D(u)$ of Theorem 8.5 has an infinite right-sided derivative at $u=1$.
(b) Generalizing the result of (a), show that if $u=\sigma(n) / n$ for some $n$, then $D(u)$ has an infinite right-sided derivative at $u$. Hint: Consider numbers of the form $n m$, where $\sigma(m) / m$ is very close to 1 .
(c) Combining part (b) with Exercise 34, prove that $D(u)$ is strictly increasing on $[1, \infty)$.
36. (Suggested by C. Pomerance) Prove that the numbers from the set $\left\{\frac{s(n)}{n}\right\}_{n \geq 2}$ have vanishing geometric mean, i.e., that $\left(\prod_{j=2}^{N} \frac{s(n)}{n}\right)^{\frac{1}{N-1}} \rightarrow 0$ as $N \rightarrow \infty$. Hint: The result of Exercise 33 may be useful.

Remark. W. Bosma \& B. Kane have considered the geometric mean of the same sequence extended only over even numbers $n$. (Note that when $n$ is even, $s(n) / n \geq 1 / 2$.) They show that this mean exists and is strictly less than 1 (in fact, it is $\approx 0.969$ ). This result, as well as the
result of the exercise, is of use in heuristic arguments surrounding the Catalan-Dickson conjecture.
37. (Erdős \& Turán [Erd45]; see also [Dre72b], [Bat72]) Let $S(x):=$ $\#\{m \in \mathbf{N}: \sigma(m) \leq x\}$.
(a) Show that $S(x) / x \rightarrow \int_{1}^{\infty} D(u) / u^{2} d u$ as $x \rightarrow \infty$.
(b) Show that the limit in (a) can also be written in the form

$$
\prod_{p}\left(1-\frac{1}{p}\right)\left(1+\frac{1}{p+1}+\frac{1}{p^{2}+p+1}+\frac{1}{p^{3}+p^{2}+p+1}+\ldots\right)
$$

Hint: Let $r(n)$ denote the number of solutions $m$ to $\sigma(m)=n$. Observe that $\sum_{n=1}^{\infty} \frac{r(n)}{n^{s}}=\sum_{m=1}^{\infty} \frac{1}{\sigma(m)^{s}}=\prod_{p}\left(\sum_{j=0}^{\infty} \frac{1}{\sigma\left(p^{j}\right)^{s}}\right)$ for real $s>1$. Apply the Dirichlet-Dedekind theorem of Exercise 7.2.
38. Given a set of natural numbers $S$, let $M(S)$ be the set of natural numbers possessing a divisor from $S$, i.e.,

$$
M(S):=\{m \in \mathbf{N}: n \mid m \text { for some } n \in S\} .
$$

For obvious reasons, $M(S)$ is referred to as the set of multiples of $S$. If $M(S)$ has an asymptotic density, we call $S$ a Besicovitch set. This (somewhat confusing) terminology honors A. S. Besicovitch, who was the first to produce, in [Bes34], an example of a set $S$ for which the asymptotic density of $M(S)$ does not exist.
(a) Show that if $S$ is finite, then $S$ is a Besicovitch set.
(b) For each $x>0$, put $S_{x}:=S \cap[1, x]$ and $S^{x}:=S \backslash S_{x}$. Show that if the upper density of $M\left(S^{x}\right)$ tends to zero as $x \rightarrow \infty$, then $S$ is Besicovitch, and in fact the density of $M(S)$ is the limit of the numbers $d_{x}$ as $x \rightarrow \infty$, where $d_{x}$ denotes the density of $M\left(S_{x}\right)$.
(c) Using the result of (b), show that if the sum of the reciprocals of the elements of $S$ converges, then $S$ is Besicovitch.
39. (Continuation; cf. Erdős [Erd70], Benkoski \& Erdős [BE74]) A natural number $n$ is said to be pseudoperfect if some subset of the proper divisors of $n$ sums to $n$. For example, 104 is pseudoperfect, since

$$
104=52+26+13+8+4+1
$$

(a) Say that the natural number $n$ is primitive pseudoperfect if $n$ is pseudoperfect and no proper divisor of $n$ is pseudoperfect, and let $S$ be the set of primitive pseudoperfect numbers. Show that the set of all pseudoperfect numbers is the set $M(S)$.
(b) Write $S=S_{1} \cup S_{2}$, where $S_{1}:=\left\{n \in S: \Omega(n)>\frac{101}{100} \log \log n\right\}$ and $S_{2}:=S \backslash S_{1}$. Using the result of Exercise 3.25, show that the upper density of $M\left(S_{1}^{x}\right)$ tends to zero as $x \rightarrow \infty$.
(c) We now turn our attention to $S_{2}$. In this part and the next, we show that the sum of the reciprocals of the elements of $S_{2}$ converges, so that the upper density of $M\left(S_{2}^{x}\right)$ tends to zero as $x \rightarrow \infty$.
For a natural number $n>1$, write $P(n)$ for the largest prime divisor of $n$. Using the result of Exercise 3.32, show that as $x \rightarrow$ $\infty$, the number of $n \leq x$ with $P(n) \leq x^{1 /(\log \log x)^{2}}$ is at most $x \exp \left(-\left(\frac{1}{2}+o(1)\right)(\log \log x)^{2}\right)$.
(d) Suppose $n \in S_{2} \cap[1, x]$ and $P(n)>x^{1 /(\log \log x)^{2}}$. Put $p=P(n)$ and write $n=p n^{\prime}$. Since $n$ is pseudoperfect, we can write $n=$ $d_{1}+d_{2}+\cdots+d_{t}$ (say), where $d_{1}, \ldots, d_{t}$ are proper divisors of $n$. By considering this decomposition modulo $p$ and using that $n$ is primitive pseudoperfect, show that $p$ divides the sum of a nonempty collection of divisors of $n^{\prime}$.
Deduce that for each fixed $n^{\prime} \leq x^{1-1 /(\log \log x)^{2}}$, the number of possibilities for $p$ is $\ll 2^{\tau\left(n^{\prime}\right)} \log x$. Using the bound $\tau\left(n^{\prime}\right) \leq 2^{\Omega\left(n^{\prime}\right)}$, deduce that the number of elements of $S_{2} \cap[1, x]$ with $P(n)>$ $x^{1 /(\log \log x)^{2}}$ is at most

$$
x \exp \left(-(1+o(1)) \log x /(\log \log x)^{2}\right) .
$$

Combining this with the result of (c), show that the sum of the reciprocals of the elements of $S_{2}$ converges.
(e) Combining (a)-(d), prove that $S$ is Besicovitch, i.e., that the set of pseudoperfect numbers possesses an asymptotic density.
40. (Benkoski \& Erdős, ibid.) It is clear that every pseudoperfect number $n$ (as defined in Exercise 39) is nondeficient, i.e., perfect or abundant. A natural number $n$ which is nondeficient but not pseudoperfect is called weird. The sequence of weird numbers begins $70,836,4030,5830, \ldots$.

Suppose that $n$ is a weird number.
(a) Show that there is no solution to

$$
1=\frac{1}{d_{1}}+\frac{1}{d_{2}}+\frac{1}{d_{3}}+\cdots+\frac{1}{d_{t}},
$$

where $d_{1}, \ldots, d_{t}>1$ are distinct divisors of $n$.
(b) Let $\epsilon$ be the smallest positive number of the form $1-\left(\frac{1}{d_{1}}+\cdots+\frac{1}{d_{t}}\right)$, with the $d_{i}$ as in (a). Show that if $m \in \mathbf{N}$ and $m n$ is not weird, then $\sigma_{-1}(m n) \geq \sigma_{-1}(n)+\epsilon$. Hint: Begin by writing $1=\sum 1 / d_{i}$, where the $d_{i}$ are distinct divisors of $m n$ and each $d_{i}>1$.
(c) Deduce from (b) and Theorem 8.5 that the set of weird multiples of $n$ has positive lower density.
Here are two open questions about weird numbers: Are all weird numbers even? Can $\sigma(n) / n$ be arbitrarily large for weird $n$ ?

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[^0]:    ${ }^{1}$ Here a Dirichlet series is a series of the form $F(s)=\sum_{n=1}^{\infty} c_{n} / n^{s}$, where each $c_{n}$ is a complex number.

[^1]:    ${ }^{2}$ Here the angle determined by $e_{1}$ and $e_{2}$ means the closed set of points between $e_{1}$ and $e_{2}$.

[^2]:    ${ }^{1}$ An arithmetic function $f$ is termed additive if $f(m n)=f(m)+f(n)$ whenever $\operatorname{gcd}(m, n)=1$.

[^3]:    ${ }^{1}$ The sole rationale for insisting that $\mathcal{A}$ be a sequence instead of a set is to ensure that duplicate elements are counted with multiplicity. Notationally, we will treat $\mathcal{A}$ below as if it were a set, but the reader should understand that $\mathcal{A}$ is actually a multiset.

[^4]:    ${ }^{1}$ One can recover that the perfect numbers make up a set of density zero from the continuity of $D(u)$, since it is plain that the upper density of the perfect numbers is bounded by $D(2+\epsilon)-D(2-\epsilon)$ for each $\epsilon>0$.

[^5]:    ${ }^{2}$ Recall from Chapter 1 that a number is said to be $y$-smooth if all of its prime factors are $\leq y$.

[^6]:    ${ }^{3}$ Recall that a divisor $m$ of $n$ is said to be unitary if $\operatorname{gcd}(m, n / m)=1$.

