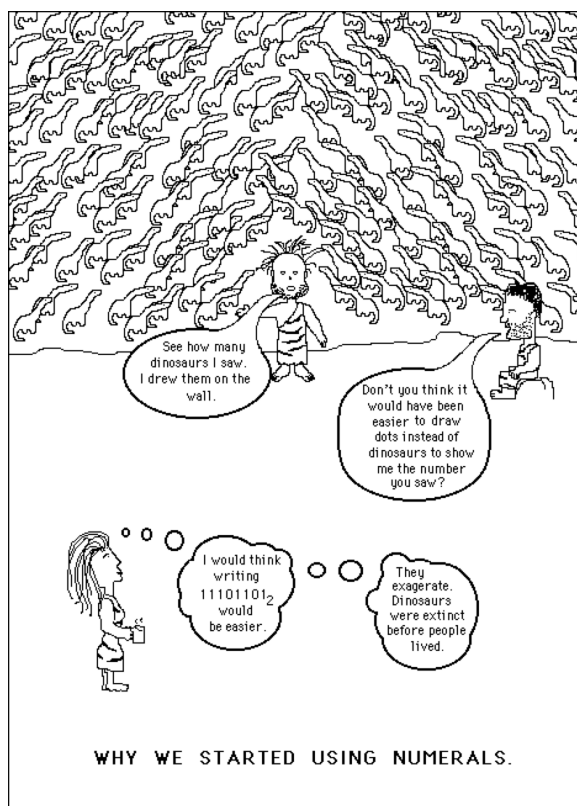# Math 788M: Computational Number Theory
(Instructor's Notes)

The Very Beginning:

- A positive integer n can be written in n steps.

- The role of numerals (now $O(\log n)$ steps)

- Can we do better? (Example: The largest known prime contains 258716 digits and doesn't take long to write down. It's $2^{859433} - 1$.)

WHY WE STARTED USING NUMERALS.

Running Time of Algorithms:

- A positive integer n in base b contains $[\log_b n] + 1$ digits.

- Big-Oh & Little-Oh Notation (as well as $\ll, \gg, \sim, \asymp$)

Examples 1: $\log(1 + (1/n)) = O(1/n)$

Examples 2: $[\log_b n] + 1 \asymp \log n$

Examples 3: $1 + 2 + \cdots + n \ll n^2$

Examples 4: f a polynomial of degree $k \implies f(n) = O(n^k)$

Examples 5: $(r+1)^\pi \sim r^\pi$

  • We will want algorithms to run quickly (in a small number of steps) in comparison to the length of the input. For example, we may ask, "How quickly can we factor a positive integer n?" One considers the length of the input n to be of order $\log n$ (corresponding to the number of binary digits n has). An algorithm runs in polynomial time if the number of steps it takes is bounded above by a polynomial in the length of the input. An algorithm to factor n in polynomial time would require that it take $O\left((\log n)^k\right)$ steps (and that it factor n).

Addition and Subtraction (of n and m):

  • We are taught how to do binary addition and subtraction in $O(\log n + \log m)$ steps.

  • We aren't going to do better than this.

  • Converting to base 10 (or any other base) is another story.

Multiplication (of n and m):

  • We are taught how to do multiplication in $O\left((\log n)(\log m)\right)$ steps.

  • Better is Possible

Example: Let $M(d)$ denote an upper bound on the number of steps required to multiply two numbers with $\leq d$ binary digits. For simplicity, we suppose n and m both have $2r$ digits. We show that nm can be calculated in $\ll r^{\log 3/\log 2} \ll r^{1.585}$ steps. Write $n = a \times 2^r + b$ and $m = c \times 2^r + d$. Then $nm = x \times 2^{2r} + y + z \times 2^r$ where $x = ac$, $y = bd$, and $z = (a+b)(c+d) - x - y$. We deduce $M(2r) \leq 3M(r+2) + kr$ for some constant k (where we have allowed for the possibility that r is not an integer). Etc.

  • Even Better is Possible

Theorem. $M(d) \ll d(\log d) \log \log d$.

  • Note that multiplying a d digit number by 19 takes $O(d)$ steps.

Sketch of proof that $M(d) \ll d^{1+\varepsilon}$.

  • **Theorem.** *Given distinct numbers $x_0, x_1, \ldots, x_k$ and numbers $y_0, y_1, \ldots, y_k$, there is a unique polynomial f of degree $\leq k$ such that $f(x_j) = y_j$ for all j.*

  • Lagrange Interpolation:

$$f(x) = \sum_{i=0}^{k} \left(\prod_{\substack{0 \leq j \leq k \\ j \neq i}} \frac{x - x_j}{x_i - x_j} \right) y_i$$

  • Mimic the Example above. Suppose n and m have $\leq kr$ digits. Write

$$n = \sum_{j=0}^{k-1} a_j 2^{jr} \quad \text{and} \quad m = \sum_{j=0}^{k-1} b_j 2^{jr}.$$

Consider

$$f(x) = \left(\sum_{j=0}^{k-1} a_j x^j \right) \left(\sum_{j=0}^{k-1} b_j x^j \right),$$

and note $nm = f(2^r)$. The coefficients of $f(x)$ can be determined by using $x_j = j$ in the Lagrange interpolation formula, taking $2k-1$ multiplications of $\leq r + c_k$ (for some constant c_k) digit numbers (to obtain the y_j). This leads to

$$M(kr) \leq (2k-1)M(r+c_k) + c_k'r \qquad \Longrightarrow \qquad M(d) \ll (2k-1)^{\log_k d} \ll d^{\log(2k-1)/\log k},$$

for some constant c_k', which implies what we want.

- What about the required division?

Homework:

(1) Let

$$f_1(n) = \log n, \quad f_2(n) = \log \log n, \quad f_3(n) = \log(3n+5),$$

$$f_4(n) = 1, \quad \text{and} \quad f_5(n) = n.$$

Determine the largest subset S of $\{1, 2, \ldots, 5\}$ for which each of (a), (b), and (c) holds (each letter below is a separate problem requiring a possibly different set S). You do not need to justify your answers.

(a) $f_1(n) = O(f_j(n))$ for $j \in S$.
(b) $f_1(n) = o(f_j(n))$ for $j \in S$.
(c) $f_1(n) \sim f_j(n)$ for $j \in S$.

(2) Find an explicit function $f(x)$ involving logarithms, powers, and/or exponentials which has both of the properties (at the same time):

(i) $f(n) \ll n^\varepsilon$ for every $\varepsilon > 0$.
(ii) $f(n) \gg (\log n)^k$ for every $k > 0$.

(Hint: I did put exponentials in the list above for a reason.) Justify your answer.

(3) The value of $f(n) = \displaystyle\sum_{k=1}^{n} \frac{1}{k}$ can be estimated by comparing it's value to an integral. For example, by comparing the sum of the areas of the rectangles indicated in the graph below with the area under the graph of $y = 1/x$, one obtains

The graph of $y = 1/x$.
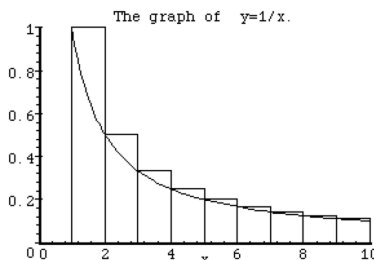
$$f(9) \geq \int_1^{10} \frac{1}{x} \, dx = \log 10.$$

(a) Prove that $f(n) \gg \log n$.

(b) Prove that $f(n) \ll \log n$. (Hint: Try another picture as above with the rectangles completely under the curve.)

(c) Prove that $f(n) \sim \log n$.

Division:

- Problem: Given two positive integers n and m, determine the quotient q and the remainder r when n is divided by m. These should be integers satisfying

$$n = mq + r \quad \text{and} \quad 0 \le r < m.$$

- Let $D(d)$ denote an upper bound on the number of steps required to obtain q and r given n and m each have $\le d$ binary digits.

- **Theorem.** *Suppose $M(d)$ has the form $df(d)$ where $f(d)$ is an increasing function of d. Then $D(d) \ll M(d)$.*

- We need only compute $1/m$ to sufficient accuracy. Suppose n and m have $\le r$ digits. If $1/m = 0.d_1 d_2...$ (in binary) with $d_1, \ldots, d_r$ known, then $n/m = (1/2^r)(n \times d_1 d_2 \ldots d_r) + \theta$ where $0 \le \theta \le 1$. If q' represents the number formed from all but the last r digits of $n \times d_1 d_2 \ldots d_r$, then $n = mq' + \theta'$ where $0 \le \theta' < 2m$. Try $q = q'$ and $q = q' + 1$.

- Preliminaries to the Proof – Newton's Method

Example: Discuss computing $1/m$. Consider $f(x) = mx - 1$ and $f(x) = m - 1/x$. Develop the idea of beginning with a good approximation x_0 to $1/m$ and obtaining successively better ones by using the recursion $x_{n+1} = 2x_n - mx_n^2$. Note that if $x_n = (1 - \varepsilon)/m$, then $x_{n+1} = (1 - \varepsilon^2)/m$. The role of multiplication already has replaced the role of division.

- Algorithm from Knuth, Vol. 2, pp. 295-6:

Algorithm R *(High-precision reciprocal)*. Let v have the binary representation $v = (0.v_1 v_2 v_3 \ldots)_2$, where $v_1 = 1$. This algorithm computes an approximation z to $1/v$, such that

$$|z - 1/v| \le 2^{-n}.$$

R1. [Initial Approximation] Set $z \leftarrow \frac{1}{4}\lfloor 32/(4v_1 + 2v_2 + v_3) \rfloor$ and $k \leftarrow 0$.

R2. [Newtonian iteration] (At this point we have a number z of the binary form $(xx.xx \ldots x)_2$ with $2^k + 1$ places after the radix point, and $z \le 2$.) Calculate $z^2 = (xxx.xx \ldots x)_2$ exactly, using a high-speed multiplication routine. Then calculate $V_k z^2$ exactly, where $V_k = (0.v_1 v_2 \ldots v_{2^{k+1}+3})_2$. Then set $z \leftarrow 2z - V_k z^2 + r$, where $0 \le r < 2^{-2^{k+1}-1}$ is added if necessary to "round up" z so that it is a multiple of $2^{-2^{k+1}-1}$. Finally, set $k \leftarrow k + 1$.

R3. [Test for end] If $2^k < n$, go back to step R2; otherwise the algorithm terminates.

- Relate algorthm to our problem (note notational differences).
- Let z_k be the value of z after k iterations of step R2. Show by induction that

$$z_k \le 2 \quad \text{and} \quad |z_k - 1/v| \le 2^{-2^k}. \tag{$*$}$$

For $k = 0$, try each possibility for v_1, v_2, and v_3. For the RHS of $(*)$, use

$$\frac{1}{v} - z_{k+1} = v\left(\frac{1}{v} - z_k\right)^2 - z_k^2\left(v - V_k\right) - r.$$

For the LHS of $(*)$, consider each of the cases $V_k = 1/2$, $V_{k-1} = 1/2 \neq V_k$, and otherwise.

- Deduce that the number of steps is

$$2M(4n) + 2M(2n) + 2M(n) + 2M(n/2) + \cdots + O(n) \ll M(n).$$

Homework:

We have sketched the argument that for every $\varepsilon > 0$, $M(d) \ll d^{1+\varepsilon}$. Most of the details given dealt with showing $M(kr) \leq (2k-1)M(r + c_k) + c_k' r$. For $k = 2$, we gave a more detailed argument. Do the same (give details) for $k = 3$ to prove that $M(d) \ll d^{\log 5/\log 3}$. Your account should explain the connection between the estimates $M(3r) \leq 5M(r + c) + c'r$ and $M(d) \ll d^{\log 5/\log 3}$. Also, clarify what the 5 multiplications are and determine a specific value for c (not c').

Elementary Number Theory:

- Modulo Arithmetic (definition, properties, & different notation)
- Computing $a^m \pmod{n}$
- Euler's Phi Function (definition, formula)
- Euler's Theorem, Fermat's Little Theorem, and Existence of Inverses
- Computing Inverses (later – see two sections from now)
- Chinese Remainder Theorem
- Generators exist modulo 2, 4, p^e, and $2p^e$

Homework:

(1) A very wealthy person buys items priced at \$2,259.29, \$4,855.07, \$9,921.23, \$12,009.91, and \$20,744.39. How many of each item did he purchase if the total purchase without taxes comes to \$749,518.05? (Hint: Factor and use a calculator or computer.)

(2) Show that for every positive integer k, there is an integer n such that each of the numbers $n+1, n+2, \ldots, n+k$ is divisible by a square > 1. For example, for $k = 3$, one can take $n = 47$ since 48 is divisible by 4, 49 is divisible by 49, and 50 is divisible by 25.

Greatest Common Divisors:

- Algorithm from Knuth, Vol. 2, p. 320:

Algorithm A *(Modern Euclidean algorithm).* Given nonnegative integers u and v, this algorithm finds their greatest common divisor. (*Note:* The greatest common divisor of *arbitrary* integers u and v may be obtained by applying this algorithm to $|u|$ and $|v|$ ....)

A1. [Check v = 0] If $v = 0$, the algorithm terminates with u as the answer.

A2. [Take u mod v] Set $r \leftarrow u \mod v$, $u \leftarrow v$, $v \leftarrow r$, and return to A1. (The operations of this step decrease the value of v, but they leave $\gcd(u,v)$ unchanged.)

- Why does the algorithm work?
- How long does it take? Explain the worse case and establish the following:

Theorem (Lamé). *Let $\phi = (1 + \sqrt{5})/2$. Let $0 \leq u, v < N$ in Algorithm A. Then the number of times step A2 is repeated is $\leq [\log_\phi(\sqrt{5}N)] - 2$.*

- The above can be used to show that the Euclidean algorithm to compute $\gcd(a,b)$ with a and b positive integers $\leq N$ has running time $\ll \log^2 N$. The running time for computing $\gcd(a,b)$ can be improved however (cf. [4, p. 428]).

Theorem. *The running time for computing the greatest common divisor of two positive integers $\leq N$ is $\ll \log N (\log \log N)^2 \log \log \log N$.*

- Inverses and a method to compute them
- **Theorem.** *Given integers a and b, not both 0, there exist integers u and v such that $au + bv = \gcd(a,b)$.*

- The average number of times step A2 is repeated is $\asymp \log N$.
- The average value of $\gcd(u,v)$ is $\asymp \log N$ but "usually" it's much smaller.
- Justification. For the first part use that

$$\sum_{d|n} \phi(d) = n \quad \text{and} \quad \sum_{d \leq N} \frac{\phi(d)}{d^2} \asymp \log N.$$

Omit the details on the second of these asymptotics. For the first, note that for each divisor d of n, the positive integers having greatest common divisor d with n are precisely the $m \leq n$ of the form kd where $1 \leq k \leq n/d$ and $\gcd(k, n/d) = 1$. This implies

$$n = \sum_{d|n} \phi(n/d) = \sum_{d|n} \phi(d).$$

The average value of $\gcd(u,v)$ is $\asymp \log N$ follows by using $\displaystyle\sum_{1 \leq n \leq N} \sum_{1 \leq m \leq N} \sum_{d|n, \, d|m} \phi(d)$. For the second part, observe that the number of pairs (n,m) with $\gcd(n,m) > z$ is bounded above by

$$\sum_{d > z} \frac{N^2}{d^2} \leq \frac{N^2}{z^2} + N^2 \int_z^\infty \frac{1}{t^2} \, dt \leq \frac{2N^2}{z}.$$

Probable Primes:

- The use of Fermat's Little Theorem
- The example $341 = 11 \times 31$ but note that $3^{340} \equiv 56 \pmod{341}$
- The example $561 = 3 \times 11 \times 17$
- Some noteworthy estimates:

$$P_2(x) \leq x^{1 - \frac{\log \log \log x}{2 \log \log x}} \quad \text{and} \quad P_2(x) \geq x^{2/7} \quad \forall x \geq x_0$$

$$\pi(x) \geq \frac{x}{\log x} = x^{1 - \frac{\log \log x}{\log x}} \quad \forall x \geq 17$$

$$P_2(2.5 \times 10^{10}) = 21853 \quad \text{and} \quad \pi(2.5 \times 10^{10}) = 1091987405$$

- Teminology: pseudoprime, probable prime, industrial grade prime, absolute pseudoprime, Carmichael number

  - The equivalence of two different definitions for absolute pseudoprimes

  - There are infinitely many absolute pseudoprimes

- Strong pseudoprimes. Suppose n is an odd composite number and write $n - 1 = 2^s m$ where m is an odd integer. Then n is a *strong pseudoprime to the base* b if either (i) $b^m \equiv 1 \pmod{n}$ or (ii) $b^{2^j m} \equiv -1 \pmod{n}$ for some $j \in [0, s-1]$.

  - Strong pseudoprimes base b are pseudoprimes base b.

  - Primes p satisfy (i) or (ii) for any b relatively prime to p.

- There are no n which are strong pseudoprimes to every base b with $1 \leq b \leq n$ and $\gcd(b, n) = 1$. To see this, assume otherwise. Note that n must be squarefree. Next, consider a prime divisor q of n, and note $n/q > 1$. Let $c \in [1, q-1]$ be such that c is not a square modulo q. Let b satisfy $b \equiv 1 \pmod{n/q}$ and $b \equiv c \pmod{q}$. Then (i) cannot hold modulo q and (ii) cannot hold modulo n/q.

  - $5^{280} \equiv 67 \pmod{561}$

- The number $3215031751 = 151 \times 751 \times 28351$ is simultaneously a strong pseudoprime to each of the bases 2, 3, 5, and 7. It's the only such number $\leq 2.5 \times 10^{10}$.

Homework:

(1) Using the Euclidean algorithm, calculate $\gcd(7046867, 1003151)$.

(2) Calculate integers x and y for which $7046867x + 1003151y = \gcd(7046867, 1003151)$.

(3) Prove that if n is an odd pseudoprime, then $2^n - 1$ is a pseudoprime.

(4) Prove that 1729 is an absolute pseudoprime.

The Lucas-Lehmer Primality Test:

- Fix integers P and Q. Let $D = P^2 - 4Q$. Define recursively u_n and v_n by

$$u_0 = 0, \quad u_1 = 1, \quad u_{n+1} = Pu_n - Qu_{n-1} \text{ for } n \geq 1,$$

$$v_0 = 2, \quad v_1 = P, \quad \text{and} \quad v_{n+1} = Pv_n - Qv_{n-1} \text{ for } n \geq 1.$$

If p is an odd prime and $p \nmid PQ$ and $D^{(p-1)/2} \equiv -1 \pmod{p}$, then $p | u_{p+1}$.

- Compute modulo p by using

$$\begin{pmatrix} u_{n+1} & v_{n+1} \\ u_n & v_n \end{pmatrix} = M^n \begin{pmatrix} 1 & P \\ 0 & 2 \end{pmatrix} \quad \text{where} \quad M = \begin{pmatrix} P & -Q \\ 1 & 0 \end{pmatrix}.$$

- The formulas

$$u_n = \frac{\alpha^n - \beta^n}{\alpha - \beta} \quad \text{and} \quad v_n = \alpha^n + \beta^n \quad \text{for } n \geq 0,$$

where $\alpha = (P + \sqrt{D})/2$ and $\beta = (P - \sqrt{D})/2$

- The formula

$$2^{n-1}u_n = \binom{n}{1}P^{n-1} + \binom{n}{3}P^{n-3}D + \binom{n}{5}P^{n-5}D^2 + \cdots$$

- Prove $p|u_{p+1}$ (let $n = p + 1$ above)

Maple's "isprime" Routine (Version 5, Release 3):

- Don't try isprime(1093^2) or isprime(3511^2) in Maple V, Release 3. The algorithm will end up in an infinite loop. This is not a concern in the latest version of Maple.

- What is isprime doing?

- The help output for isprime:

FUNCTION: isprime - primality test

CALLING SEQUENCE:
    isprime(n)

PARAMETERS:
    n - integer

SYNOPSIS:

- The function isprime is a probabilistic primality testing routine.

- It returns false if n is shown to be composite within within one strong pseudo-primality test and one Lucas test and returns true otherwise. If isprime returns true, n is "very probably" prime - see Knuth "The art of computer programming", Vol 2, 2nd edition, Section 4.5.4, Algorithm P for a reference and H. Reisel, "Prime numbers and computer methods for factorization". No counter example is known and it has been conjectured that such a counter example must be hundreds of digits long.

SEE ALSO: nextprime, prevprime, ithprime

- The Maple program:

```
proc (n)
    local btor, nr, p, r;
    options remember, system,
        'Copyright 1993 by Waterloo Maple Software';
    if not type(n,integer) then
        if type(n,numeric) then
            ERROR('argument must be an integer')
        else
            RETURN('isprime(n)')
        fi
    fi;
    if n < 2 then
        false
    elif has('isprime/w',n) then
        true
    elif igcd(2305567963945518424753102147331756070,n) <> 1 then
```

```
            false
        elif n < 10201 then
            true
        elif igcd(849696948923341811053233990918734996592606258664893273366
            11545426342203893270769390909069477309509137509786917118668028861499333382
            50976823867229837379629630667576741311267365789364407881571869698937306633
            11306647862044862494925732402262739543736363903875260816675866125595683463
            30697220447512298848222285500626837863425199602259963013159456444700647206
            9662175047724452891592786711,n) <> 1 then
            false
        elif n < 1018081 then
            true
        else nr := igcd(408410100000,n-1);
            nr := igcd(nrˆ5,n-1);
            r := iquo(n-1,nr);
            btor := modp(power(2,r),n);
                if 'isprime/cyclotest'(n,btor,2,r) = false
                    or irem(nr,3) = 0 and 'isprime/cyclotest'(n,btor,3,r) = false
                    or irem(nr,5) = 0 and 'isprime/cyclotest'(n,btor,5,r) = false
                    or irem(nr,7) = 0 and 'isprime/cyclotest'(n,btor,7,r) = false then
                    RETURN(false)
                fi;
            for p from 3 while (numtheory[jacobi])(pˆ2-4,n) <> -1 do od;
            evalb('isprime/TraceModQF'(p,n+1,n) = [2, p])
        fi
end
```

Some of the Subroutines Used:

   • I found this reference years later to some of the subroutines mentioned above. I only "suspect" that these are very close to the programs that were used in the version of "isprime" above.

```
cyclotest := proc(n,btor,i,r)
local nr, k, t, s, x;
option
    'Copyright (c) 1993 Gaston Gonnet, Wissenschaftliches Rechnen, ETH Zurich. All rights reserved.';
nr := iquo(n-1,r);
for k from 0 while irem(nr,i,'t') = 0 do nr := t od;
x := modp(power(btor,nr), n);
if x=1 then RETURN(FAIL) fi;
to k do
    s := sumxtor(x,i,n);
    if s[1]=0 then RETURN(FAIL) elif s[2]=1 then RETURN(false) fi;
    x := s[2]
od;
false
```

end:

```
sumxtor := proc(x,r,n)
    local i1, i2, i3, r1, r2, r3, s, xj;
    option
'Copyright (c) 1993 Gaston Gonnet, Wissenschaftliches Rechnen, ETH Zurich. All rights reserved.';
    if r < 1 then ERROR('r is too small')
    elif r = 1 then [1,x]
    elif r < 6 then
       s := 1+x;
       xj := modp(x^2,n);
       to r-2 do s := s+xj; xj := modp(xj*x,n) od;
       [modp(s,n),xj]
    else i1 := isqrt(r);
       i2 := iquo(r,i1,'i3');
       if i3=0 then
          r2 := procname(x,i2,n);
          r1 := procname(r2[2],i1,n);
          [modp(r2[1]*r1[1],n), r1[2]]
       else r3 := procname(x,i3,n);
          r2 := procname(x,i2,n);
          r1 := procname(r2[2],i1,n);
          [modp( r3[1] + r3[2] * modp(r2[1]*r1[1],n), n), modp(r3[2]*r1[2],n) ]
       fi;
    fi;
end:

TraceModQF := proc ( p, k, n )
local i, kk, trc, v;
option
    'Copyright (c) 1993 Gaston Gonnet, Wissenschaftl. Rechnen, ETH Zurich. All rights reserved.';
kk := k;
for i while kk > 1 do kk := iquo(kk+1,2,v[i]) od;
trc := [ p, 2 ];
for i from i-1 by -1 to 1 do
    if v[i]=1 then
       trc := modp( [trc[1]^2 - 2, trc[1]*trc[2] - p], n );
    else trc := modp( [trc[1]*trc[2] - p, trc[2]^2 - 2], n );
    fi
od;
trc
end:
```

Maple's Initial Steps:

- The list 'isprime/w' consists of the primes < 100.

- The number $2305567\ldots6070$ is the product of the primes < 100. The next prime is 101 and $101^2 = 10201$.

- The number $8496969\ldots7113$ is the product of the primes in the interval $(100, 1000)$. The next prime is 1009 and $1009^2 = 1018081$.

- If n is a prime < 1018081, the initial steps will declare it to be prime. If n is a composite number with a prime factor < 1000, then the initial steps will declare it composite. In particular, all $n < 1018081$ will have been properly dealt with.

Maple's Version of the Strong-Pseudoprime Test (or is it?):

- $408410100000 = 2^5 3^5 5^5 7^5$

- If $n - 1 = 2^{e_1} 3^{e_2} 5^{e_3} 7^{e_4} m$ where $\gcd(210, m) = 1$, then $nr = 2^{e'_1} 3^{e'_2} 5^{e'_3} 7^{e'_4}$ where $e'_j = \min\{e_j, 25\}$. Also, $r = (n-1)/nr$ and $btor = 2^r \mod n$.

- What is 'isprime/cyclotest'(n,btor,2,r) doing? It computes $x = 2^y \mod n$ where $y = (n-1)/2^{e'_1}$. If $x \equiv \pm 1 \pmod{n}$, then it tells the algorithm to go on to the next step (involving 'isprime/cyclotest'(n,btor,3,r)). Otherwise, it replaces x with $x^2 \mod n$. If now $x = 1$, then it declares n composite. If $x \equiv -1 \pmod{n}$, then it tells the algorithm to go on to the next step. If $x \not\equiv \pm 1 \pmod{n}$, then it again replaces x with $x^2 \mod n$. It continues like this, declaring n composite if $x = 1$, telling the algorithm to go on if $x \equiv -1 \pmod{n}$, and replacing x with $x^2 \mod n$ otherwise, until e'_1 squarings of x have occured. At that point, if the algorithm has not declared n to be a prime or composite number and it has not told the algorithm to go on, then it declares that n is composite. This is equivalent to a strong pseudoprime test to the base 2 if $e'_1 = e_1$ (for most numbers this will be the case). Otherwise it is a little weaker.

- The steps involving 'isprime/cyclotest'(n,btor,i,r) for $i \in \{3, 5, 7\}$ are variations on a strong pseudoprime test to the base 2 (yes, I mean base 2). For example, 'isprime/cyclotest'(n,btor,3,r) makes use of the fact that if $x^3 \equiv 1 \pmod{p}$, then either $x \equiv 1 \pmod{p}$ or $x^2 + x + 1 \equiv 0 \pmod{p}$ (this is analogous to using for the strong pseudoprime test that if $x^2 \equiv 1 \pmod{p}$, then either $x \equiv 1 \pmod{p}$ or $x \equiv -1 \pmod{p}$).

- The "help" for isprime is somewhat misleading. The references cited do not suggest (as far as I noticed - OK, I didn't read every word in these references) that one should use e'_1 above rather than e_1. Algorithm P of Knuth's book is not used (it involves choosing an integer $b \in (1, n)$ at random and checking if n is a strong pseudoprime base b).

Maple's Version of the Lucas-Lehmer Test:

- Take $Q = 1$. Then

$$v_{2n} = v_n^2 - 2 \quad \text{and} \quad v_{2n+1} = v_{n+1} v_n - P \quad \text{for } n \geq 1.$$

Also, $Du_n = 2v_{n+1} - Pv_n$. If p is a prime, then $v_p \equiv P \pmod{p}$. Isprime checks if (v_{n+1}, v_n) is congruent to $(2, P)$ modulo n.

- Prove that if it is, then $n | u_{n+1}$.

- How is (v_{n+1}, v_n) computed modulo n? Beginning with $\vec{w} = (v_1, v_0)$ and considering the binary digits of n beginning from the left-most digit, $\vec{w} = (v_{m+1}, v_m)$ is replaced by (v_{2m+2}, v_{2m+1}) whenever the digit 1 is encountered and by (v_{2m+1}, v_{2m}) whenever the digit 0 is encountered. In this way, the subscript of the second coordinate of $\vec{w}$ corresponds to the number obtained by considering a left portion of the binary representation of n.

The Fix to this "isprime" Routine:

- The error occurs due to the possibility that n is a perfect square. Before the Lucas-Lehmer test at the end of the program, one can simply add the line:

  if isqrt(n)^2=n then RETURN(false) fi;

Mersenne Primes:

- Definition.
- It's connection to perfect numbers.
- **The Lucas-Lehmer Test.** *Let p be an odd prime, and define recursively*

$$L_0 = 4 \quad and \quad L_{n+1} = L_n^2 - 2 \mod (2^p - 1) \ for \ n \geq 0.$$

Then $2^p - 1$ is a prime if and only if $L_{p-2} = 0$.

- The Lucas sequence with $Q = 1$ and $P = 4$. Here, $L_n = v_{2^n} \mod (2^p - 1)$. Note that

$$u_n = \frac{(2 + \sqrt{3})^n - (2 - \sqrt{3})^n}{\sqrt{12}} \quad and \quad v_n = (2 + \sqrt{3})^n + (2 - \sqrt{3})^n.$$

- ($\Longrightarrow$) Suppose $N = 2^p - 1$ is a prime. We will use that (from the theory of quadratic reciprocity) $3^{(N-1)/2} \equiv -1 \pmod{N}$ (where here N prime and $N \equiv 7 \pmod{12}$ is important). Also, $2^p \equiv 1 \pmod{N}$ easily implies there is an x such that $x^2 \equiv 2 \pmod{N}$. Hence, $2^{(N-1)/2} \equiv x^{N-1} \equiv 1 \pmod{N}$. We want $v_{(N+1)/4} \equiv 0 \pmod{N}$. From $v_{2n} = v_n^2 - 2$, it follows that we need only show $v_{(N+1)/2} \equiv -2 \pmod{N}$. Observe that $2 \pm \sqrt{3} = ((\sqrt{2} \pm \sqrt{6})/2)^2$. Hence,

$$v_{(N+1)/2} = \left(\frac{\sqrt{2} + \sqrt{6}}{2} \right)^{N+1} + \left(\frac{\sqrt{2} - \sqrt{6}}{2} \right)^{N+1}$$

$$= 2^{-N} \sum_{j=0}^{(N+1)/2} \binom{N+1}{2j} \sqrt{2}^{N+1-2j} \sqrt{6}^{2j} = 2^{(1-N)/2} \sum_{j=0}^{(N+1)/2} \binom{N+1}{2j} 3^j.$$

Using

$$\binom{N+1}{2j} = \binom{N}{2j} + \binom{N}{2j-1},$$

we deduce

$$2^{(N-1)/2} v_{(N+1)/2} \equiv 1 + 3^{(N+1)/2} \equiv -2 \pmod{N},$$

and the result follows.

- ($\Longleftarrow$, the more important case) We make use of the following two identities:

$$v_n = u_{n+1} - u_{n-1} \quad and \quad u_{m+n} = u_m u_{n+1} - u_{m-1} u_n, \tag{1}$$

where the subscripts are all assumed to be non-negative integers. Establish

$$if \ u_n \equiv 0 \pmod{p^e}, \ then \ u_{pn} \equiv 0 \pmod{p^{e+1}}, \tag{2}$$

where e is assumed to be a positive integer. To obtain (2), use induction to show that if $a = u_{n+1}$, then

$$u_{kn} \equiv k a^{k-1} u_n \pmod{p^{e+1}} \quad \text{and} \quad u_{kn+1} \equiv a^k \pmod{p^{e+1}}$$

and then take $k = p$; for example, observe that for $k = 2$, we have

$$u_{2n} = u_n u_{n+1} - u_{n-1} u_n = u_n u_{n+1} + u_n(u_{n+1} - 4u_n) \equiv 2 a u_n \pmod{p^{e+1}}$$

and $u_{2n+1} = u_{n+1}^2 - u_n^2 \equiv a^2 \pmod{p^{e+1}}$.

Next, we use that

$$u_n = \sum_{k=0}^{\lfloor (n-1)/2 \rfloor} \binom{n}{2k+1} 2^{n-2k-1} 3^k \quad \text{and} \quad v_n = \sum_{k=0}^{\lfloor n/2 \rfloor} \binom{n}{2k} 2^{n-2k+1} 3^k$$

to obtain that

$$u_p \equiv 3^{(p-1)/2} \pmod{p} \quad \text{and} \quad v_p \equiv 4 \pmod{p}. \tag{3}$$

Fermat's Little Theorem implies for $p > 3$ that $u_p \equiv \pm 1 \pmod{p}$. Using (1), (3), and the definition of the u_n, we obtain that if $u_p \equiv 1 \pmod{p}$, then $u_{p-1} \equiv 4u_p - u_{p+1} \equiv 4u_p - v_p - u_{p-1} \equiv -u_{p-1} \pmod{p}$ implying $u_{p-1} \equiv 0 \pmod{p}$. Similarly, if $u_p \equiv -1 \pmod{p}$, then $u_{p+1} \equiv 4u_p - u_{p-1} \equiv 4u_p + v_p - u_{p+1} \equiv -u_{p+1} \pmod{p}$ implying $u_{p+1} \equiv 0 \pmod{p}$. Thus, for every prime $p > 3$, there is an integer $\epsilon = \epsilon(p) = \pm 1$ such that

$$u_{p+\epsilon} \equiv 0 \pmod{p}. \tag{4}$$

Also, $u_p \equiv 0 \pmod{p}$ if $p = 2$ and $p = 3$.

Observe that

$$\gcd(u_n, u_{n+1}) = 1 \quad \text{and} \quad \gcd(u_n, v_n) \le 2. \tag{5}$$

The former follows from the recursive definition of u_n. The second follows from the first by first noting $2u_{n+1} = 4u_n + v_n$ (obtained by combining (1) with the recursive relation on u_n).

For any positive integer m, denote by $\alpha = \alpha(m)$ the smallest positive integer for which $u_\alpha \equiv 0 \pmod{m}$ (it is not needed, but such an α always exists). Then

$$u_n \equiv 0 \pmod{m} \quad \Longleftrightarrow \quad \alpha | n. \tag{6}$$

This follows by considering $u_\alpha, u_{\alpha+1}, \ldots$ modulo m.

By assumption, $v_{2^{p-2}} \equiv L_{p-2} \equiv 0 \pmod{2^p - 1}$. Thus, (5) $\implies u_{2^{p-2}} \not\equiv 0 \pmod{2^p - 1}$. Also, $u_{2n} = u_n v_n$ implies $u_{2^{p-1}} \equiv 0 \pmod{2^p - 1}$. It follows that $\alpha(2^p - 1) = 2^{p-1}$.

Write $2^p - 1 = p_1^{e_1} p_2^{e_2} \cdots p_r^{e_r}$ with p_j distinct primes and e_j positive integers. Each $p_j \ge 3$. Set

$$k = \operatorname{lcm}\left\{ p_j^{e_j - 1}(p_j + \epsilon_j) : j = 1, \ldots, r \right\}.$$

Here $\epsilon_j = \pm 1$ are chosen so that (4) holds with $p = p_j$ and $\epsilon = \epsilon_j$. Observe that (2), (4), and (6) imply $u_k \equiv 0 \pmod{2^p - 1}$. In particular, k is a multiple of $\alpha(2^p - 1) = 2^{p-1}$. By the definition of k, it follows that 2^{p-1} divides $p_j^{e_j - 1}(p_j + \epsilon_j)$ for some j. For such j, we have $p_j \ge 2^{p-1} - 1$. The inequality

$$3p_j \ge 2^p + 2^{p-1} - 3 > 2^p - 1$$

now implies that $2^p - 1$ must be prime.

Homework:

(1) Prove that each of the identities in (1) holds for arbitrary positive integers n and m.

(2) Prove that for every positive integer m, the number $\alpha = \alpha(m)$ in the proof above exists.

General Primality Tests:

- A polynomial time algorithm may exist:

Theorem (Selfridge-Weinberger). *Assume the Extended Riemann Hypothesis holds. Let n be an odd integer > 1. A necessary and sufficient condition for n to be prime is that for all positive integers $a < \min\{70(\log n)^2, n\}$, we have $a^{(n-1)/2} \equiv \pm 1 \pmod{n}$ with at least one occurrence of -1.*

- For $n = 1729$ and all integers $a \in [1, n-1]$ with $\gcd(a, n) = \gcd(a, 7 \times 13 \times 19) = 1$, $a^{(n-1)/2} \equiv 1 \pmod{n}$.

- Observe that $a < n$ exists when n is prime (consider a primitive root).

- **Theorem (Lucas).** *Let n be a positive integer. If there is an integer a such that $a^{n-1} \equiv 1 \pmod{n}$ and for all primes p dividing $n - 1$ we have $a^{(n-1)/p} \not\equiv 1 \pmod{n}$, then n is prime.*

- For primes such an a exists (consider a primitive root).

- Prove the theorem. Note that one can obtain that n is a prime if for each p dividing $n - 1$, there is an integer a (possibly depending on p) such that $a^{n-1} \equiv 1 \pmod{n}$ and $a^{(n-1)/p} \not\equiv 1 \pmod{n}$.

- **Theorem (Pepin Test).** *Let $F_n = 2^{2^n} + 1$ with n a positive integer. Then F_n is prime if and only if $3^{(F_n - 1)/2} \equiv -1 \pmod{F_n}$.*

- Prove the theorem.

- **Theorem (Proth, Pocklington, Lehmer Test).** *Let n be a positive integer. Suppose $n - 1 = FR$ where all the prime factors of F are known and $\gcd(F, R) = 1$. Suppose further that there exists an integer a such that $a^{n-1} \equiv 1 \pmod{n}$ and for all primes p dividing F we have $\gcd(a^{(n-1)/p} - 1, n) = 1$. Then every prime factor of n is congruent to 1 modulo F.*

- Consider the case $F \geq \sqrt{n}$.

- Prove the theorem. (Let q be a prime divisor of n. Consider $m = \text{ord}_q(a)$. Show $F|m$ by showing $p^e||F$ and $p^e \nmid m$ is impossible.)

- In 1980, Adleman, Pomerance, and Rumely developed a primality test that determines if n is prime in about $(\log n)^{c \log \log \log n}$ steps (shown by Odlyzko).

Factoring Algorithms (Part I):

- Given a composite integer $n > 1$, the general problem is to find some nontrivial factorization of n, say $n = uv$ where each of u and v is an integer > 1. If this can be done effectively and one has a good primality test, one will have a good method for completely factoring n.

- The expectation is that a random number n will have on the order of $\log \log n$ prime factors. Describe what this means but don't prove it.

- Most numbers n have a prime factor $> \sqrt{n}$. Prove using

$$\sum_{p \leq x} \frac{1}{p} = \log \log x + A + O(1/\log x).$$

- One expects typically small prime factors, so it is reasonable to first do a quick "sieve" to determine if this is the case.

- **Pollard's $p - 1$ Factoring Algorithm**

This algorithm determines a factorization of a number n if n has a prime factor p where $p - 1$ factors into a product of small primes. A simple form of the algorthm is to compute $2^{k!} \mod n$ successively for $k = 1, 2, \dots$ until some prescribed amount (say 10^6 or 10^7), for each k checking $\gcd(2^{k!} - 1 \mod n, n)$ to possibly obtain a nontrivial factor of n. If $(p - 1)|k!$, then $2^{k!} \equiv 1 \pmod{p}$ so that p will divide $\gcd(2^{k!} - 1 \mod n, n)$ and the chances of obtaining a nontrivial factor of n will be good.

- **Pollard's ρ-Algorithm**

This method typically finds a prime factor p of n in about $\sqrt{p}$ steps (so $O(N^{1/4})$ steps). Note that small prime factors will usually be found first.

- Preliminary observation. Suppose we roll a fair die with n faces k times. We claim that if $k \geq 2\sqrt{n} + 2$, then with probability $> 1/2$ two of the numbers rolled will be the same. (Mention the birthday problem.) Use that there are at least $\sqrt{n}$ integers in $[\sqrt{n}, 2\sqrt{n} + 1]$. The probability the numbers rolled are all different is

$$\prod_{j=1}^{k-1} \left(\frac{n - j}{n} \right) \leq \left(1 - \frac{\sqrt{n}}{n} \right)^{\sqrt{n}} \leq \frac{1}{e}.$$

So the result follows.

- The algorithm. Let $f(x) = x^2 + 1$. Let $f^{(j)}$ be defined by $f^{(1)}(x) = f(x)$ and $f^{(j+1)}(x) = f(f^{(j)}(x))$ for $j \geq 1$. Compute $a_j = f^{(j)}(1) \mod n$ for $1 \leq j \leq k$. The idea is that if $k \geq 100\sqrt{p}$ where p is a prime factor of n, then one can expect to find i and j with $1 \leq i < j \leq k$ such that $f^{(i)}(1) \mod p = f^{(j)}(1) \mod p$. In this case, we would have $a_i \equiv a_j \pmod{p}$ so that p will divide $\gcd(a_i - a_j, n)$. Hopefully then by computing $\gcd(a_i - a_j, n)$ we can determine a factorization of n. This will not be so good however if we actually compute $\binom{k}{2}$ different gcd's.

We instead use Floyd's cycle-finding algorithm. Observe that if i and j are as above, then (since there are $j - i$ integers in $(i, j]$) there is an integer $t \in (i, j]$ for which $(j - i)|t$. Now, $a_i \equiv a_j \pmod{p}$ implies $a_{i+u} \equiv a_{j+u} \pmod{p}$ for every positive integer u. In particular,

$$a_t \equiv a_{t+(j-i)} \equiv a_{t+2(j-i)} \equiv a_{t+3(j-i)} \equiv \cdots \equiv a_{2t} \pmod{p}.$$

Thus, rather than checking $\binom{k}{2}$ different gcd's as above, one can compute $a_1, a_2, \dots$ and check as one progresses the values of $\gcd(a_{2t} - a_t, n)$ for $t = 1, 2, \dots$. One continues until one finds a factorization of n, noting again that this should take $O(\sqrt{p})$ steps to find a given prime factor p.

- Brent and Pollard factored $F_8 = 2^{2^8} + 1$ using this method with $f(x) = x^{1024} + 1$. Discuss why such a choice for $f(x)$ would be appropriate here.

Dixon's Factoring Algorithm:

- The basic idea. Suppose $n = p_1^{e_1} p_2^{e_2} \cdots p_r^{e_r}$ with p_j odd distinct primes and $e_j \in \mathbb{Z}^+$. Then $x^2 \equiv 1 \pmod{p_j^{e_j}}$ has two solutions implies $x^2 \equiv 1 \pmod{n}$ has 2^r solutions. If x and y are random and $x^2 \equiv y^2 \pmod{n}$, then with probability $(2^r - 2)/2^r$ we can factor n (nontrivially) by considering $\gcd(x + y, n)$.

- The algorithm.

(1) Randomly choose a number a and compute $s(a) = a^2 \mod n$. (We want $a > \sqrt{n}$.)

(2) A bound $B = B(n)$ to be specified momentarily is chosen. Determine if $s(a)$ has a prime factor $> B$. We choose a new a if it does. Otherwise, we obtain a complete factorization of $s(a)$.

(3) Let $p_1, \ldots, p_t$ denote the primes $\leq B$. We continue steps (1) and (2) until we obtain $t + 1$ different a's, say $a_1, \ldots, a_{t+1}$.

(4) From the above, we have the factorizations
$$s(a_i) = p_1^{e(i,1)} p_2^{e(i,2)} \cdots p_t^{e(i,t)} \quad \text{for } i \in \{1, 2, \ldots, t + 1\}.$$
Compute the vectors
$$\vec{v}_i = \langle e(i, 1), e(i, 2), \ldots, e(i, t) \rangle \mod 2 \quad \text{for } i \in \{1, 2, \ldots, t + 1\}.$$
These vectors are linearly dependent modulo 2. Use Gaussian elimination (or something better) to find a non-empty set $S \subseteq \{1, 2, \ldots, t + 1\}$ such that $\sum_{i \in S} \vec{v}_i \equiv \vec{0} \pmod{2}$. Calculate $x \in [0, n - 1] \cap \mathbb{Z}$ (in an obvious way) satisfying
$$\prod_{i \in S} s(a_i) \equiv x^2 \pmod{n}.$$

(5) Calculate $y = \prod_{i \in S} a_i \mod n$. Then $x^2 \equiv y^2 \pmod{n}$. Compute $\gcd(x + y, n)$. Hopefully, a nontrivial factorization of n results.

- A (non-realistic) example. Take $n = 1189$ and $B = 11$. Suppose after step (3) we have for the a_i's the numbers 151, 907, 449, 642, 120, and 1108 with the $s(a_i)$'s being $210 = 2 \times 3 \times 5 \times 7$, $1050 = 2 \times 3 \times 5^2 \times 7$, $660 = 2^2 \times 3 \times 5 \times 11$, $770 = 2 \times 5 \times 7 \times 11$, $132 = 2^2 \times 3 \times 11$, and $1108 = 2^3 \times 7 \times 11$, respectively. Execute the rest of the algorithm.

Homework:

Use Dixon's Algorithm to factor $n = 80099$. Suppose $B = 15$ and the a_j's from the first three steps are the numbers 1392, 58360, 27258, 39429, 12556, 42032, and 1234. (Each of these squared reduced modulo n should have all of its prime factors $\leq B$.)

- What about B? To analyze the optimal choice for B, we use
$$\psi(x, y) = |\{n \leq x : p|n \implies p \leq y\}|.$$
From previous discussions, $\psi(x, \sqrt{x}) \sim (1 - \log 2)x$. In general, $\psi(x, x^{1/u}) \sim \rho(u)x$ for some number $\rho(u)$.

Theorem 1 (Dickman). *For u fixed, $\psi(x, x^{1/u}) \sim \rho(u)x$ where $\rho(u)$ satisfies:*
    (i) $\rho(u)$ is continuous for $u > 0$
    (ii) $\rho(u) \to 0$ as $u \to \infty$
    (iii) $\rho(u) = 1$ for $0 < u \leq 1$
    (iv) for $u > 1$, $\rho(u)$ satisfies the differential delay equation $u\rho'(u) = -\rho(u-1)$.

Helmut Maier, improving on work of deBruijn, essentially removes any restriction on u by establishing the result above whenever $u < (\log x)^{1-\varepsilon}$ for any fixed $\varepsilon > 0$. Note that $x^{1/\log x} = e$. Also, deBruijn established that

$$\rho(u) = \exp\left(-(1 + o(1))u \log u\right) \approx \frac{1}{u^u}.$$

We are interested in $\psi(n, B)$. Thus, $u = \log n / \log B$. Note that $u \leq \log n$. Suppose $u < (\log n)^{1-\varepsilon}$ as above. We deduce

$$\psi(n, B) = n \exp\left(-(1 + o(1)) \log n \log u / \log B\right).$$

The number of different times we expect to go through steps (1) and (2) of the algorithm is

$$(\pi(B) + 1) \exp\left((1 + o(1)) \log n \log u / \log B\right).$$

We expect $\leq B$ steps to factor each value of $s(a)$. We are led to considering

$$B = \exp\left(\sqrt{\log n}\sqrt{\log u}/\sqrt{2}\right) = \exp\left(\sqrt{\log n}\sqrt{\log \log n}/2\right).$$

The number of steps expected for Dixon's Algorithm is therefore $\exp\left(2\sqrt{\log n}\sqrt{\log \log n}\right)$ (take into account the Gaussian elimination).

- Preliminaries to the CFRAC Algorithm. Discuss simple continued fractions (scf). Mention that if a/b is a reduced convergent of the scf for α, then $|\alpha - (a/b)| < 1/b^2$. Deduce that $|a^2 - \alpha^2 b^2| \ll \alpha$.

- The CFRAC Algorithm. Use Dixon's Algorithm with the a_j's chosen so that a_j is the numerator of a reduced convergent of the scf for $\sqrt{n}$. If b_j is the denominator, then $|a_j^2 - nb_j^2| < 2\sqrt{n}$. In other words, if one modifies $s(a_j)$ so that it might be negative, we can take $|s(a_j)| < 2\sqrt{n}$. One can deal with negative $s(a_j)$ by treating -1 as a prime. The running time is improved as $\psi(n, B)$ above gets to be replaced by $\psi(2\sqrt{n}, B)$. One obtains here a running time of $O\left(\exp(\sqrt{2}\sqrt{\log n}\sqrt{\log \log n})\right)$. On the other hand, it has not been established that the values of $s(a_j)$ here are just as likely to have all prime factors $\leq B$ as random numbers their size, so this running time is heuristic. We note that this algorithm was used by Brillhart and Morrison to factor $F_7 = 2^{2^7} + 1$.

- An "early abort" strategy can be combined with the above ideas to reduce the running time of the algorithms. Given a, one stops trying to factor $s(a)$ if it has no "small" prime factors. This leads to a running time of the form $O\left(\exp(\sqrt{3/2}\sqrt{\log n}\sqrt{\log \log n})\right)$.

Homework:

Use the CFRAC Algorithm to factor $n = 135683$. The first 10 numerators of the convergents for the simple continued fraction for $\sqrt{n}$ and their squares modulo n (in factored form) are shown below. You may use this information. As far as B goes, ignore it (take it to be 1000 if you want). You do not need to use Gaussian elimination. Instead you can simply try to find the right combination of factors mentally.

Numerators	Squares Mod n
737	19×23
1105	$-1 \times 2 \times 61$
6262	257
13629	$-1 \times 2 \times 193$
19891	11×23
53411	$-1 \times 2 \times 7 \times 11$
233535	277
520481	$-1 \times 2 \times 7 \times 19$
1274497	11×19
4343972	-1×83

- The Quadratic Sieve Algorithm. Consider $F(x) = (x + [\sqrt{n}])^2 - n$. Then $|F(x)| \ll |x|\sqrt{n}$ for $0 < |x| \leq \sqrt{n}$. The idea of the quadratic sieve algorithm is to consider a in Dixon's Algorithm to be of the form $a = x + [\sqrt{n}]$ with $|x|$ small. Here, we allow for $s(a)$ to be negative as in the CFRAC Algorithm. Thus, $|s(a)| \ll |x|\sqrt{n}$.

- Why would this be better than the CFRAC Algorithm? In the CFRAC Algorithm, we had $|s(a)| < 2\sqrt{n}$, so this is a reasonable question. The advantage of the Quadratic Sieve Algorithm is that one can "sieve" prime divisors of $s(a)$ to determine how $s(a)$ factors for many a at once. To clarify, for a small prime p, one can solve the quadratic $F(x) \equiv 0 \pmod{p}$. If there are solutions to the congruence, there will usually be two incongruent solutions modulo p, say x_1 and x_2. Thus, one knows that if $x \equiv x_1$ or x_2 modulo p and $a = x + [\sqrt{n}]$, then $p|s(a)$. Otherwise, $p \nmid s(a)$.

- The running time is $O\left(\exp(\sqrt{9/8}\sqrt{\log n}\sqrt{\log \log n})\right)$ for the Quadratic Sieve Algorithm. The running time is heuristic.

- There are other variations of the above algorithms. In particular, a version of the Quadratic Sieve Algorithm suggested by Peter Montgomery reduces the running time to

$$O\left((1/2)\sqrt{\log n}\sqrt{\log \log n}\right).$$

The Number Field Sieve:

- Let f be an irreducible monic polynomial with integer coefficients. Let α be a root of f. Let m be an integer for which $f(m) \equiv 0 \pmod{n}$. The mapping $\phi : \mathbb{Z}[\alpha] \to \mathbb{Z}_n$ with $\phi(g(\alpha)) = g(m) \mod n$ for all $g(x) \in \mathbb{Z}[x]$ is a homomorphism. The idea is to find a set S of polynomials $g(x) \in \mathbb{Z}[x]$ such that (i) $\prod_{g \in S} g(m) = y^2$ for some $y \in \mathbb{Z}$, and (ii) $\prod_{g \in S} g(\alpha) = \beta^2$

for some $\beta \in \mathbb{Z}[\alpha]$. Taking $x = \phi(\beta)$, we deduce

$$x^2 \equiv \phi(\beta)^2 \equiv \phi(\beta^2) \equiv \phi\left(\prod_{g \in S} g(\alpha)\right) \equiv \prod_{g \in S} g(m) \equiv y^2 \pmod{n}.$$

Thus, we can hope once again to factor n by computing $\gcd(x + y, n)$.

• What do we choose for f? We determine an f of degree $d > 1$ and with $n > 2^{d^2}$ as follows. Set $m = [n^{1/d}]$. Write n in base m; in other words, compute $c_0, c_1, \ldots, c_d$ each in $\{0, 1, \ldots, m-1\}$ with

$$n = c_d m^d + c_{d-1} m^{d-1} + \cdots + c_1 m + c_0.$$

Set $f(x) = \sum_{j=0}^{d} c_j x^j$. Then f is monic and $f(m) \equiv 0 \pmod{n}$. Next, we attempt to factor f in $\mathbb{Z}[x]$. (Discuss the running time for this.) If it is irreducible, then we are happy. If it is reducible, then we use the factorization of f to determine a factorization of n (explain).

Homework:

Prove that f is monic. In other words, show that $c_d = 1$.

• What do we choose for the g? We take the g to be of the form $a - bx$ where $|a| \le D$ and $0 < b \le D$. We want $g(m)$ to have only small prime factors. This is done by first choosing b and then, with b fixed, letting a vary and sieving to determine the a for which $g(m)$ has only small prime factors.

• How do we obtain the desired square in $\mathbb{Z}[\alpha]$? Let $\alpha_1, \ldots, \alpha_d$ be the distinct roots of f with $\alpha = \alpha_1$. We consider the norm map $\mathrm{N}(g(\alpha)) = g(\alpha_1) \cdots g(\alpha_d)$, where $g \in \mathbb{Z}[x]$. It has the two properties: **(i)** if g and h are in $\mathbb{Z}[x]$, then $\mathrm{N}(g(\alpha)h(\alpha)) = \mathrm{N}(g(\alpha))\,\mathrm{N}(h(\alpha))$, and **(ii)** if $g \in \mathbb{Z}[x]$, then $\mathrm{N}(g(\alpha)) \in \mathbb{Z}$. It follows that the norm of a square in $\mathbb{Z}[\alpha]$ is a square in $\mathbb{Z}$. On the other hand,

$$\mathrm{N}(a - b\alpha) = b^d \prod_{j=1}^{d} \left(\frac{a}{b} - \alpha_j\right) = b^d f(a/b) = a^d + c_{d-1}a^{d-1}b + \cdots + c_1 a b^{d-1} + c_0 b^d.$$

The idea is to try to obtain a set S of pairs (a, b) as above. As we force the product $\prod(a - bm)$ to be a square (products over $(a, b) \in S$), we also force $\prod\left(a^d + c_{d-1}a^{d-1}b + \cdots + c_0 b^d\right)$ to be a square.

• Have we really answered the previous question? No. We have found an element of $\mathbb{Z}[\alpha]$ with its norm being a square in $\mathbb{Z}$. This does not mean the element is a square in $\mathbb{Z}[\alpha]$. (For example, consider $\alpha = i$ and $S = \{(2, 1), (2, -1)\}$.) There is quite a bit of work left in the algorithm in terms of modifying the exponent vectors to produce the squares we want, but a small bit of the ideas have been demonstrated.

• The running time for the number field sieve is $\exp\left(c(\log n)^{1/3}(\log\log n)^{2/3}\right)$ where c is a constant and $c = 4/(3^{2/3})$ will do.

• The number field sieve was used to factor $F_9 = 2^{2^9} + 1$ in 1993 by Lenstra, Lenstra, Manasse, and Pollard.

Public-Key Encryption:

- **Problem:** *Describe how to communicate with someone you have never met before through the personals without anyone else understanding the private material you are sharing with this stranger.*

- What they don't know can hurt them. Find two large primes p and q; right now 100 digit primes will suffice. Compute $n = pq$. If you are secretive about your choices for p and q, then you can tell the world what n is and you will know something no one else in the world knows, namely how n factors. Now, don't you feel special?

- What to publish first. Now that you know something no one else knows, the problem becomes a bit easier to resolve. You (and you alone) can determine $\phi(n)$. Do so, and choose some positive integer s (the "encrypting exponent") with $\gcd(s, \phi(n)) = 1$. You publish n and s in the personals. You also describe to someone how to correspond with you in such a way that no one else will understand the message (at this point everyone can read what you are writing, but that's OK as it is only a temporary situation). What you tell them is something like this: to form a message M concatenate the symbols 00 for blank, 01 for a, 02 for b, ..., 26 for z, 27 for a comma, 28 for a period, and whatever else you might want. For example, $M = 0805121215$ means "hello". Next, tell the person to publish (back in the personals) the value of $E = M^s \mod n$. (The person should be told to make sure that $M^s > n$ by adding extra blanks if necessary and that $M < n$ by breaking up a message into two or more messages if necessary.)

- What can you do with the encoded message? You calculate t with $st \equiv 1 \pmod{\phi(n)}$ (one can use $t \equiv s^{\phi(\phi(n))-1} \mod \phi(n)$). Then compute $E^t \mod n$. This will be the same as M modulo n (unless p or q divides M, which isn't likely). So now you know the message.

- Computing $\phi(n)$ is seemingly as difficult as factoring n. Here one needs to compute $\phi(\phi(n))$ which for you is only as difficult as factoring $p - 1$ and $q - 1$. If these each have 100 digits, then n has around 200 digits and all is reasonable (for the times we live in). One can also try to construct p and q with $p - 1$ and $q - 1$ of some nice form (for example, having small prime factors).

Homework:

Let $p = 193$ and $q = 257$, so $n = pq = 49601$. Let $s = 247$.
(a) Someone sends the encrypted message $E = 48791$. Determine the word sent.
(b) Encrypt the message, "No". In other words, tell me the value of E.

- Certified signatures. Imagine person A has published n and s in the personals, person B is corresponding with person A in private, and person C really dislikes person B. C decides to send A a message in the personals that reads something like, "Dear A, I think you are a jerk. Your dear friend, B." This of course would make A very upset with B and would make C very happy. What would be nice is if there were a way for B to sign his messages so that A can see the signature and know whether a message supposedly from B is really from B. This is done as follows. First, if B is really corresponding with A, he (B stands for boy) should have his very own n and s which he has shared with at least A. Let's call them n' and s', and let the corresponding t be t'. Now, B doesn't have to use his name, but he informs A of some signature (name) he will use, say S. He can change S regularly if he wishes, but in any case, it is given to A as part of an encrypted

message. At the end of the encrypted message, he gives A the number $T = S^{t'} \mod n'$. After A decodes the message, he computes $T^{s'} \mod n'$ (remember n' and s' are public). The result will be S. Since only B knows t', only B could have determined T and A will know that the message really came from B.

Factoring Polynomials:

- **Problem:** *Given a polynomial in $\mathbb{Z}[x]$, determine if it is irreducible over $\mathbb{Q}$. If it is reducible, find a non-trivial factorization of it in $\mathbb{Z}[x]$.*

- **Berlekamp's Algorithm.** This algorithm determines the factorization of a polynomial $f(x)$ in $\mathbb{Z}_p[x]$ where p is a prime (or more generally over finite fields). For simplicity, we suppose $f(x)$ is monic and squarefree in $\mathbb{Z}_p[x]$. Let $n = \deg f(x)$. Let A be the matrix with jth column derived from the coefficients reduced modulo p of $x^{(j-1)p} \mod f(x)$. Specifically, write

$$x^{(j-1)p} \equiv \sum_{i=1}^{n} a_{ij} x^{i-1} \mod f(x) \qquad \text{for } 1 \le j \le n.$$

Then we set $A = (a_{ij} \mod p)_{n \times n}$. Note that the first column consists of a one followed by $n-1$ zeroes. In particular, $\langle 1, 0, 0, \ldots, 0 \rangle$ will be an eigenvector for A associated with the eigenvalue 1. We are interested in determining the complete set of eigenvectors associated with the eigenvalue 1. In other words, we would like to know the null space of $B = A - I$ where I represents the $n \times n$ identity matrix. It will be spanned by $k = n - \text{rank}(B)$ linearly independent vectors which can be determined by performing row operations on B. Suppose $\vec{v} = \langle b_1, b_2, \ldots, b_n \rangle$ is one of these vectors, and set $g(x) = \sum_{j=1}^{n} b_j x^{j-1}$. Observe that $g(x^p) \equiv g(x) \pmod{f(x)}$ in $\mathbb{Z}_p[x]$. Moreover, the $g(x)$ with this property are precisely the $g(x)$ with coefficients obtained from the components of vectors $\vec{v}$ in the null space of B.

Claim. $f(x) \equiv \prod_{s=0}^{p-1} \gcd(g(x) - s, f(x)) \pmod{p}$.

The proof follows by using that $f(x)$ divides $g(x)^p - g(x) \equiv \prod_{s=0}^{p-1}(g(x) - s)$ in $\mathbb{Z}_p[x]$ and that each irreducible factor of $f(x)$ in $\mathbb{Z}_p[x]$ divides at most one of the $g(x) - s$. Observe that if $g(x)$ is not a constant, then $1 \le \deg(g(x) - s) < \deg f(x)$ for each s so the above claim implies we get a non-trivial factorization of $f(x)$ in $\mathbb{Z}_p[x]$. On the other hand, $f(x)$ will not necessarily be completely factored. One can completely factor $f(x)$ by repeating the above procedure for each factor obtained from the claim; but it is simpler to use (and not difficult to show) that if one takes the product of the greatest common divisors of each factor of $f(x)$ obtained above with $h(x) - s$ (with $0 \le s \le p-1$) where $h(x)$ is obtained from another of the k vectors spanning the null space of B, then one will obtain a new non-trivial factor of $f(x)$ in $\mathbb{Z}_p[x]$. Continuing to use all k vectors will produce a complete factorization of $f(x)$ in $\mathbb{Z}_p[x]$. (As an example of Berlekamp's algorithm, factor $x^7 + x^4 + x^3 + x + 1$ in $\mathbb{Z}_2[x]$.)

- **Hensel Lifting Algorithm.** This algorithm gives a method for using the factorization of $f(x)$ in $\mathbb{Z}_p[x]$ (p a prime) to produce a factorization of $f(x)$ in $\mathbb{Z}_{p^k}[x]$. Suppose that $u(x)$ and $v(x)$ are relatively prime polynomials in $\mathbb{Z}_p[x]$ for which

$$f(x) \equiv u(x)v(x) \pmod{p}.$$

Then Hensel Lifting will produce for any positive integer k polynomials $u_k(x)$ and $v_k(x)$ in $\mathbb{Z}[x]$ satisfying

$$u_k(x) \equiv u(x) \pmod{p}, \quad v_k(x) \equiv v(x) \pmod{p},$$

and

$$f(x) \equiv u_k(x)v_k(x) \pmod{p^k}.$$

When $k = 1$, it is clear how to choose $u_k(x)$ and $v_k(x)$. For $k \geq 1$, we determine values of $u_{k+1}(x)$ and $v_{k+1}(x)$ from the values of $u_k(x)$ and $v_k(x)$ as follows. We compute

$$w_k(x) \equiv \frac{1}{p^k}\left(f(x) - u_k(x)v_k(x)\right) \pmod{p}.$$

Since $u(x)$ and $v(x)$ are relatively prime in $\mathbb{Z}_p[x]$, we can find $a(x)$ and $b(x)$ in $\mathbb{Z}_p[x]$ (depending on k) such that

$$a(x)u(x) + b(x)v(x) \equiv w_k(x) \pmod{p}.$$

It follows that we can take

$$u_{k+1}(x) = u_k(x) + b(x)p^k \quad \text{and} \quad v_{k+1}(x) = v_k(x) + a(x)p^k.$$

A complete factorization of $f(x)$ modulo p^k can be obtained from a complete factorization of $f(x)$ modulo p by modifying this idea. However, note that $f(x) = 2x^2 + 5x + 3 = (x+1)(2x+3)$ satisfies

$$f(x) \equiv \left(x + \frac{3^k + 3}{2}\right)(2x + 2) \pmod{3^k}.$$

Homework:

(1) Use Berlekamp's algorithm to factor $f(x) = x^6 + x^3 + x^2 + x + 1$ modulo 2. You should obtain two polynomials $u(x)$ and $v(x)$ of degrees < 6 such that $f(x) \equiv u(x)v(x) \pmod{2}$.

(2) Use the previous problem and Hensel lifting to factor $f(x) = x^6 + x^3 + x^2 + x + 1$ modulo 32. To help, let $u(x)$ and $v(x)$ be as in the previous problem with $\deg u > \deg v$. Then you can take advantage of the following:

$$u(x) + x^3v(x) \equiv x^5 + 1 \pmod{2}$$

$$(x+1)u(x) + (x^2+1)v(x) \equiv x^5 + x^4 \pmod{2}$$

$$xu(x) + (x^2+x)v(x) \equiv x^5 \pmod{2}.$$

• **An Inequality of Landau.** This inequality gives an upper bound on the "size" of the factors of a given polynomial in $\mathbb{Z}[x]$. Given $f(x) = \sum_{j=0}^{n} a_j \in \mathbb{Z}[x]$, we measure it's size with $\|f\| = \left(\sum_{j=0}^{n} a_j^2\right)^{1/2}$. Thus for a fixed $f(x) \in \mathbb{Z}[x]$, we want an upper bound on $\|g\|$ where $g(x)$ is a factor of $f(x)$ in $\mathbb{Z}[x]$.

Theorem. *If $f(x)$, $g(x)$, and $h(x)$ in $\mathbb{Z}[x]$ are such that $f(x) = g(x)h(x)$, then $\|g\| \leq 2^{\deg g}\|f\|$.*

For $f(x) = a_n \prod_{j=1}^{n}(x - \alpha_j)$ with $a_n \neq 0$ and α_j complex but not necessarily distinct, we define the Mahler measure of f as $M(f) = |a_n| \prod_{j=1}^{n} \max\{1, |\alpha_j|\}$. The following two properties are easily seen to hold: (i) if $g(x)$ and $h(x)$ are in $\mathbb{C}[x]$, then $M(gh) = M(g)M(h)$; (ii) if $g(x)$ is in $\mathbb{Z}[x]$, then $M(g) \geq 1$.

Claim. For $f(x) \in \mathbb{R}[x]$, $M(f) \leq \|f\| \leq 2^{\deg f} M(f)$.

For the claim, we define for a given $w(x) \in \mathbb{C}[x]$, $\tilde{w}(x) = x^{\deg w} w(1/x)$. The coefficient of $x^{\deg w}$ in the expanded product $w(x)\tilde{w}(x)$ is $\|w\|^2$. For $f(x) = \sum_{j=0}^{n} a_j x^j = a_n \prod_{j=1}^{n}(x - \alpha_j)$, we consider

$$w(x) = a_n \prod_{\substack{1 \leq j \leq n \\ |\alpha_j| > 1}} (x - \alpha_j) \prod_{\substack{1 \leq j \leq n \\ |\alpha_j| \leq 1}} (\alpha_j x - 1).$$

One checks that

$$w(x)\tilde{w}(x) = a_n^2 \prod_{j=1}^{n}(x - \alpha_j) \prod_{j=1}^{n}(1 - \alpha_j x) = f(x)\tilde{f}(x).$$

By comparing coefficients of x^n, we deduce that $\|w\| = \|f\|$. Also, observe that from the definition of w, $|w(0)| = M(f)$. Thus, if $w = \sum_{j=0}^{n} c_j x^j$, then

$$M(f) = |c_0| \leq (c_0^2 + c_1^2 + \cdots + c_n^2)^{1/2} = \|w\| = \|f\|,$$

establishing the first inequality. For the second inequality observe that for any $k \in \{1, 2, \ldots, n\}$, the product of any k of the α_j has absolute value $\leq M(f)/|a_n|$. It follows that $|a_{n-k}|/|a_n|$, which is the sum of the products of the roots taken k at a time, is $\leq \binom{n}{k} \times M(f)/|a_n|$. Hence, $|a_{n-k}| \leq \binom{n}{k} M(f) = \binom{n}{n-k} M(f)$. The second inequality now follows from

$$\|f\| = \left(\sum_{j=0}^{n} a_j^2\right)^{1/2} \leq \sum_{j=0}^{n} |a_j| \leq \sum_{j=0}^{n} \binom{n}{j} M(f) = 2^n M(f).$$

To prove the theorem, just use the Claim and properties (i) and (ii) of Mahler measure to deduce

$$\|g\| \leq 2^{\deg g} M(g) \leq 2^{\deg g} M(g) M(h) = 2^{\deg g} M(gh) = 2^{\deg g} M(f) \leq 2^{\deg g} \|f\|.$$

- **Combining the above ideas.** We factor a given $f(x) \in \mathbb{Z}[x]$ with the added assumptions that it is monic and squarefree. The latter we can test by computing $\gcd(f, f')$, which will give us a nontrivial factor of f if f is not squarefree. If f were not monic, a little more needs to be added to the ideas below (but not much).

Let $B = 2^{(\deg f)/2} \|f\|$. Then if f has a nontrivial factor g in $\mathbb{Z}[x]$, it has such a factor of degree $\leq (\deg f)/2$ so that by Landau's inequality, we can use B as a bound on $\|g\|$.

Next, we find a prime p for which f is squarefree modulo p. There are a variety of ways this can be done. There are only a finite number of primes for which f is not squarefree modulo p (these primes divide the resultant of f and f'). Working with $\gcd(f, f')$ modulo p can resolve the issue or simply using Berlekamp's factoring algorithm until a squarefree factorization occurs is fine.

We choose a positive integer r as small as possible such that $p^r > 2B$. One factors f modulo p by Berlekamp's algorithm and uses Hensel lifting to obtain the factorization of f modulo p^r. Given our conditions on f, we can suppose all irreducible factors are monic and do so.

Now, we can determine if $f(x) = g(x)h(x)$ for some monic g and h in $\mathbb{Z}[x]$ with $\|g\| \leq B$ as follows. We observe that the coefficients of g are in $[-B, B]$. We use a residue system modulo p^r that includes this interval, namely $(-p^r/2, p^r/2]$, and consider each factorization of f modulo

p^r with coefficients in this residue system as a product of two monic polynomials $u(x)$ and $v(x)$. Since $f = gh$, there must be some factorization where $g \equiv u \pmod{p^r}$ and $h \equiv v \pmod{p^r}$. On the other hand, the coefficients of g and u are all in $(-p^r/2, p^r/2]$ so that the coefficients of $g - u$ are each divisible by p^r and are each $< p^r$ in absolute value. This implies $g = u$. Thus, we can determine if a factor g exists as above by simply checking each monic factor of f modulo p^r with coefficients in $(-p^r/2, p^r/2]$.

- Swinnerton-Dyer's Example. The algorithm just described above for factoring a polynomial $f(x) \in \mathbb{Z}[x]$ of degree n can take time that is exponential in n for some, albeit rare, $f(x)$. This has been illustrated by a nice example due to Swinnerton-Dyer. We formulate this way. Let $a_1, a_2, \ldots, a_m$ be arbitrary squarefree pairwise relatively prime integers > 1. Let S_m be the set of 2^m different m-tuples $(\varepsilon_1, \ldots, \varepsilon_m)$ where each $\varepsilon_j \in \{1, -1\}$. Then the polynomial

$$f(x) = \prod_{(\varepsilon_1, \ldots, \varepsilon_m) \in S_m} \left(x - (\varepsilon_1 \sqrt{a_1} + \cdots + \varepsilon_m \sqrt{a_m}) \right)$$

has the properties:

(i) The polynomial $f(x)$ is in $\mathbb{Z}[x]$.

(ii) It is irreducible over the rationals.

(iii) It factors into a product of linear and quadratic polynomials modulo every prime p.

To see (i), use elementary symmetric functions. Also, note that if $f_m(x)$ is the polynomial $f(x)$ above corresponding to a list $a_1, a_2, \ldots$, then $f_m(x) = f_{m-1}(x + \sqrt{a_m}) f_{m-1}(x - \sqrt{a_m})$.

To see (ii), use induction on m. The irreducibility of $f(x)$ can be deduced from showing that $\mathbb{Q}(\sqrt{a_1}, \ldots, \sqrt{a_m})$ is an extension of degree 2^m over $\mathbb{Q}$ whose Galois group over $\mathbb{Q}$ is generated by the automorphisms $\sigma_1, \ldots, \sigma_m$ defined by $\sigma_j(\sqrt{a_j}) = -\sqrt{a_j}$ and $\sigma_j(\sqrt{a_i}) = \sqrt{a_i}$ for $i \neq j$. The induction is completed by showing that

$$\sqrt{a_{m+1}} \notin \mathbb{Q}(\sqrt{a_1}, \ldots, \sqrt{a_m}).$$

Assume otherwise. Then $\sqrt{a_{m+1}}$ can be expressed as a linear combination over $\mathbb{Q}$ of the 2^m numbers formed by taking products of the numbers $\sqrt{a_j}$ (where $1 \leq j \leq m$). Apply σ_1 to such an expression and note that the result must be $\sigma_1(\sqrt{a_{m+1}})$ which is either $\sqrt{a_{m+1}}$ or $-\sqrt{a_{m+1}}$. In either case (by multiplying by -1 in the second situation), we get "another" expression for $\sqrt{a_{m+1}}$ as a linear combination over $\mathbb{Q}$ of products of the numbers $\sqrt{a_j}$ (with $1 \leq j \leq m$). In the first case, we deduce that

$$\sqrt{a_{m+1}} \in \mathbb{Q}(\sqrt{a_2}, \ldots, \sqrt{a_m})$$

and can apply a suitable induction hypothesis. In the second case, we deduce

$$\sqrt{a_1 a_{m+1}} \in \mathbb{Q}(\sqrt{a_2}, \ldots, \sqrt{a_m})$$

and can again apply the induction hypothesis.

To see (iii), again use induction on m. Fix p. If some a_j is a square modulo p, then let $g(x)$ be the polynomial corresponding to $f(x)$ above but with the $m-1$ numbers $a_1, \ldots, a_{j-1}, a_{j+1}, \ldots, a_m$. Then $f(x) = g(x + \sqrt{a_j}) g(x - \sqrt{a_j})$. It follows that if $b^2 \equiv a_j \pmod{p}$, then $f(x) \equiv$

$g(x + b)g(x - b)$ (mod p). Since each of $g(x + b)$ and $g(x - b)$ factors as a product of linears and quadratics modulo p by the induction hypothesis, we are through in this case. Now, suppose no a_j is a square modulo p. Fix $(\varepsilon_1, \ldots, \varepsilon_m) \in S_m$ and observe that when the product

$$\left(x + (\varepsilon_1 \sqrt{a_1} + \cdots + \varepsilon_m \sqrt{a_m}) \right) \left(x - (\varepsilon_1 \sqrt{a_1} + \cdots + \varepsilon_m \sqrt{a_m}) \right)$$

is expanded the result is an expression with each radicand the product of two of the a_j. Since no a_j is a square modulo p, each such product will be (the product of two non-quadratic residues is a quadratic residue). This means that the above product can be expressed as a quadratic modulo p with coefficients from $\{0, 1, \ldots, p - 1\}$. Pairing then the linear factors of $f(x)$ appropriately leads to the desired factorization modulo p.

Discrete and Fast Fourier Transforms:

- Goal: Let $M(d)$ denote the number of binary bit operations needed to multiply two positive integers each consisting of $\leq d$ bits. We have stated without proof that $M(d) \ll d(\log d) \log \log d$. We give here the basic idea behind the proof, not worrying so much about the running time but concentrating on the main idea of using fast Fourier transforms for performing multiplication.

- A not-so-discrete formula. The following example is vaguely related to what we are after. Let R be a rectangle, and suppose R is expressed as a union of rectangles R_j, $1 \leq j \leq r$, with edges parallel to R and common points only along these edges. Suppose further that each R_j has at least one edge of integer length. Then show that R itself has an edge of integer length by using that

$$\left| \iint_R e^{2\pi i(x+y)} dx\, dy \right| = \left| \sum_{j=1}^{r} \iint_{R_j} e^{2\pi i(x+y)} dx\, dy \right|.$$

For connections to what we really want, observe that

$$\int_0^1 e^{i2\pi k\theta} d\theta = \frac{1}{2\pi} \int_0^{2\pi} e^{ik\theta} d\theta = \begin{cases} 0 & \text{if } k \neq 0 \\ 1 & \text{if } k = 0, \end{cases}$$

where k is an integer.

- Prove the following more discrete version of the above formula.

Lemma. *Let n and k be integers with $n \geq 1$. Let $\omega = e^{2\pi i/n}$. Then*

$$\sum_{j=0}^{n-1} \omega^{kj} = \begin{cases} 0 & \text{if } k \not\equiv 0 \pmod{n} \\ n & \text{if } k \equiv 0 \pmod{n}. \end{cases}$$

- Example. Show that

$$1 - \frac{1}{4} + \frac{1}{7} - \frac{1}{10} + \cdots = \frac{3\log 2 + \sqrt{3}\pi}{9}.$$

Call the sum S. Use that

$$z^2 \log(1 + z) = z^3 - \frac{z^4}{2} + \frac{z^5}{3} - \frac{z^6}{4} + \cdots .$$

Let $\omega = e^{2\pi i/3}$, and note that

$$(1 + \omega)(1 + \omega^2) = 1 + \omega + \omega^2 + \omega^3 = \omega^3 = 1.$$

This can be deduced also from

$$\omega = \frac{-1 + \sqrt{3}\,i}{2} \quad \text{and} \quad \omega^2 = \frac{-1 - \sqrt{3}\,i}{2}.$$

The lemma implies that

$$3S = \log(1 + 1) + \omega^2 \log(1 + \omega) + \omega \log(1 + \omega^2)$$

$$= \log 2 - \frac{1}{2} \log\left((1 + \omega)(1 + \omega^2)\right) + \frac{\sqrt{3}\,i}{2} \log\left((1 + \omega^2)/(1 + \omega)\right)$$

$$= \log 2 + \sqrt{3}\,i \log\left(1 + \omega^2\right) = \log 2 + \sqrt{3}i \log\left(-\omega\right) = \log 2 + \sqrt{3}\,\pi/3,$$

from which the value of S above follows. Note that this example illustrates precisely what we are after with the above lemma, a way to isolate certain terms. We used third roots of unity to isolate every third term of a series. We will similarly use nth roots of unity to isolate each term of a polynomial of degree $n - 1$.

- Definitions and notations. For n a positive integer, we set $\omega = \omega_n = e^{2\pi i/n}$. Let

$$D = D(n, \omega) = \begin{pmatrix} 1 & 1 & 1 & \cdots & 1 \\ 1 & \omega & \omega^2 & \cdots & \omega^{n-1} \\ 1 & \omega^2 & \omega^4 & \cdots & \omega^{2(n-1)} \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ 1 & \omega^{n-1} & \omega^{2(n-1)} & \cdots & \omega^{(n-1)^2} \end{pmatrix}.$$

Given a vector $\vec{u} = \langle u_0, u_1, \ldots, u_{n-1} \rangle \in \mathbb{C}^n$, we define $\vec{v} = \langle v_0, v_1, \ldots, v_{n-1} \rangle$, called the discrete Fourier transform of $\vec{u}$, as $D\,\vec{u}^T$ (the product of the matrix D with the matrix consisting of a single column with the entries of $\vec{u}$). The inverse discrete Fourier transform of a vector $\vec{v} \in \mathbb{C}^n$ is defined as $(1/n)D(n, \omega^{-1})\,\vec{v}^T$. We justify that this is truly an inverse (not just in name).

- Prove that $(1/n)D(n, \omega^{-1})D(n, \omega)$ is the identity matrix by using the lemma above.
- The polynomial connection. Observe that if $f(x) = \sum_{j=1}^{n-1} a_j x^j \in \mathbb{C}[x]$, then

$$D \langle a_0, a_1, \ldots, a_{n-1} \rangle^T = \langle f(1), f(\omega), \ldots, f(\omega^{n-1}) \rangle^T. \tag{7}$$

On the other hand, if we know the values $f(1), f(\omega), \ldots, f(\omega^{n-1})$ of a polynomial $f(x)$ and do not know the coefficients of $f(x)$, then we can obtain the coefficients from

$$D^{-1} \langle f(1), f(\omega), \ldots, f(\omega^{n-1}) \rangle^T = \langle a_0, a_1, \ldots, a_{n-1} \rangle^T. \tag{8}$$

Note that here $D^{-1} = (1/n)D(n, \omega^{-1})$.

The operations D and D^{-1} are more related than might appear. Define

$$F(x) = f(1) + f(\omega)x + \cdots + f(w^{n-1})x^{n-1}.$$

Then (8) is equivalent to the assertion that

$$F(1) = na_0, \ F(\omega^{-1}) = na_1, \ F(\omega^{-2}) = na_2, \ \ldots, F(\omega^{-(n-1)}) = na_{n-1}.$$

In other words, (8) is simply asserting that if the matrix $D(n, \omega^{-1})$ is multiplied by the coefficient vector for $F(x)$, then one gets the vector $\langle F(1), F(\omega^{-1}), \ldots, F(\omega^{-(n-1)}) \rangle^T$. Hence, (8) is essentially (7) for a different polynomial and with the role of ω replaced by ω^{-1}.

- The Fast Fourier Transform. We explain a fast way of performing the computation in (7). There exist unique polynomials f_e and f_o in $\mathbb{C}[x]$ such that $f(x) = f_e(x^2) + x f_o(x^2)$. Compute $f_e(\omega^{2j})$, $f_o(\omega^{2j})$ and $\omega^j f_o(\omega^{2j})$ for $0 \le j \le n - 1$. Use this to get the right side of (7).

- Why is this fast? Note that computing the right side of (7) in an obvious way takes n^2 multiplications and $n(n-1)$ additions. The above gives a considerable improvement over this "obvious" computation (assuming that we have nth roots of unity available to us). For convenience, view n as a power of 2. Let $A(n)$ be the number of arithmetic operations of addition and multiplication needed for the computation in (7). So $A(n)$ is the number of such operations needed to compute $f(1), f(\omega), \ldots, f(\omega^{n-1})$. Observe that

$$f_e(\omega_n^{2j}) = f_e(\omega_{n/2}^j) = f_e(\omega_{n/2}^{j+(n/2)}) = f_e(\omega_n^{2(j+(n/2))}) \qquad \text{for } 0 \le j \le (n/2) - 1.$$

Similar equations hold with f_e replaced by f_o. We deduce that computing the values of $f_e(\omega^{2j})$ and $f_o(\omega^{2j})$ takes $2A(n/2)$ arithmetic operations. Multiplying $f_o(\omega^{2j})$ by ω^j and then adding $f_e(\omega^{2j})$ takes $2n$ more arithmetic operations. We obtain

$$A(n) = 2A(n/2) + 2n \implies A(n) = (2/\log 2) \, n \log n + Cn$$

for some appropriate constant C (depending on say $A(2)$). A similar analysis holds for the number of arithmetic operations needed for the computation in (8).

- Multiplying two polynomials. An obvious way of obtaining the coefficients of a product of two polynomials of degree $\le n$ takes on the order of n^2 multiplications. This can be done considerably faster using the above ideas (and assuming, as before, that we have nth roots of unity available to us). Let

$$g(x) = b_{r-1}x^{r-1} + \cdots + b_1 x + b_0 \quad \text{and} \quad h(x) = c_{s-1}x^{s-1} + \cdots + c_1 x + c_0$$

be polynomials in $\mathbb{C}[x]$ with $s \le r$. Let $n = 2r$, and rewrite $g(x)$ and $h(x)$ with leading zeroes as

$$g(x) = b_{n-1}x^{n-1} + \cdots + b_1 x + b_0 \quad \text{and} \quad h(x) = c_{n-1}x^{n-1} + \cdots + c_1 x + c_0.$$

Let $f(x) = g(x)h(x)$, and note that $\deg f \le n - 1$. From (7), we obtain

$$D \langle b_0, b_1, \ldots, b_{n-1} \rangle^T = \langle g(1), g(\omega), \ldots, g(\omega^{n-1}) \rangle^T$$

and

$$D \langle c_0, c_1, \ldots, c_{n-1} \rangle^T = \langle h(1), h(\omega), \ldots, h(\omega^{n-1}) \rangle^T.$$

These can be computed with $2A(n)$ arithmetic operations where $A(n)$ is as before. Computing $f(\omega^j) = g(\omega^j)h(\omega^j)$ for $0 \le j \le n - 1$ takes n multiplications. Now, applying (8) allows us to

compute the coefficients of $f(x)$ with $A(n)$ further arithmetic operations. Rescaling, we deduce that the product of two polynomials of degree $\leq n$ in $\mathbb{C}[x]$ can be computed with $O(n \log n)$ complex additions and multiplications.

- Integer multiplication? Suppose we wish to multiply two positive integers A and B. Write each in binary as

$$A = (b_{r-1} \ldots b_1 b_0)_2 \quad \text{and} \quad B = (c_{s-1} \ldots c_1 c_0)_2.$$

The idea then is to take $g(x)$ and $h(x)$ as above, and compute the product $f(x) = g(x)h(x)$. Observe that the coefficients of $f(x)$ are not bits. However, the product AB can still be obtained by computing $f(2)$ which, after computing the coefficients of $f(x)$, amounts to some shifts and additions that do not require much time. The main difficulty is in having to deal with roots of unity. In 1968, Volker Straussen showed that it is possible to approximate the roots of unity to sufficient accuracy to obtain a total complexity for the multiplication of order $n \log n (\log \log n)^{1+\varepsilon}$ (for any fixed $\varepsilon > 0$) where n is a bound on the number of bits of A and B. So this is close to what is the best known bound. We will explain an alternative approach that avoids computations of the complex roots of unity.

- The main idea is to replace complex roots of unity with roots of 1 modulo some prime (or various primes) p. More precisely, let d be a positive integer (for computational purposes, it is best that d consist of only of small prime divisors; in particular, a power of 2 can be used for d). Let p be a prime such that $d | (p-1)$. Then there is a positive integer ω such that ω has order d modulo $p - 1$. Observe that if p divides a Fermat number $2^{2^k} + 1$, for example, and $d = 2^{k+1}$, then one can take $\omega = 2$. The following lemma replaces our previous one.

Lemma. *Let t be an integer. With the notation above,*

$$\sum_{j=0}^{d-1} \omega^{tj} \equiv \begin{cases} 0 \pmod{p} & \text{if } t \not\equiv 0 \pmod{d} \\ d \pmod{p} & \text{if } t \equiv 0 \pmod{d}. \end{cases}$$

- Work modulo p. Observe that d has an inverse modulo p since $p > d \geq 1$. Also, ω has an inverse, namely ω^{d-1}, modulo p. We thus define the discrete Fourier transform modulo p using ω as our root of unity and note that both $D(d, \omega)$ and its inverse $D(d, \omega^{-1})/d$ are defined modulo p. The previous arguments all carry through in this setting.

- Some analysis of the complexity (and clarification of the algorithm). Suppose that we want to multiply two positive integers A and B and construct, as before,

$$g(x) = b_{n-1} x^{n-1} + \cdots + b_1 x + b_0 \quad \text{and} \quad h(x) = c_{n-1} x^{n-1} + \cdots + c_1 x + c_0,$$

where the b_j and c_j are either 0 or 1 (though, in practice, larger coefficients can be considered), $A = g(2)$, $B = h(2)$ and the leading coefficients are stacked to be 0's in such a way that if

$$r - 1 = \max\{j : b_j \neq 0\} \quad \text{and} \quad s - 1 = \max\{j : c_j \neq 0\},$$

then $s \leq r \leq n/2$. Note that $n/2$ is a bound for the number of bits making up A and B. We consider $n = d$ and a prime p as above with $d | (p-1)$. We also let, as before, ω denote an element of order d modulo p. (Discuss how this can be found.) We compute the coefficients of $f(x) =$

$g(x)h(x)$ modulo p by using the fast Fourier transform to obtain the product of $D(d, \omega)$ with the coefficient vectors for $g(x)$ and $h(x)$ giving $g(\omega^j)$ and $h(\omega^j)$ (that is, modulo p) for $0 \leq j \leq n-1$, by computing $f(\omega^j) = g(\omega^j)h(\omega^j)$ (modulo p) for each such j, and then by using the fast Fourier transform to obtain the product of $D(d, \omega^{-1})/d$ with the vector $\langle f(1), f(\omega), \ldots, f(\omega^{n-1}) \rangle^T$.

Observe that once we know the coefficients of $f(x)$ modulo p, we can compute $f(2)$ modulo p, but what we are really wanting is the value of $f(2)$ not modulo p. We get around this by making sure p is sufficiently large (or by using similar information obtained from various primes and combining the information with an application of the Chinese Remainder Theorem). More precisely, we observe that the coefficients of $g(x)h(x)$ are each ≥ 0 and

$$\leq s \leq d/2 < p/2.$$

Hence, if the coefficients of $f(x)$ modulo p are taken to be in the interval $[0, p)$, then they will be the coefficients of $f(x)$ (not modulo p).

To examine the complexity, observe that arithmetic is done modulo p (this is replacing the need to compute approximations to complex roots of unity). So the running time analysis must take this arithmetic into consideration. We want a prime p chosen so that $d|(p-1)$. Keeping in mind that d can be chosen to be on the order of the number of bits of A and B, distribution results for primes imply that such primes can be found that are not very large (so, in any case, arithmetic modulo p can be done in time that is polynomial in the logarithm of the number of bits of A and B). The size of p can be reduced somewhat by decreasing the size of $n = d$ and increasing the size of the coefficients a_j and b_j so that $A = g(2^\ell)$ and $B = h(2^\ell)$ for some appropriate positive integer ℓ. Then the size of the coefficients of $g(x)h(x)$ may exceed p, but we can work around this by computing the coefficients of $g(x)h(x)$ modulo several primes and then using the Chinese Remainder Theorem to obtain the actual coefficients of $g(x)h(x)$.

We note the example given by John Pollard [14]. Suppose we want to multiply two 10000 bit numbers A and B. Take $\ell = 21$ so that $A = g(2^{21})$ and $B = h(2^{21})$ correspond to writing A and B in base 2^{21} (or breaking up their binary representations into blocks of 21 bits). Then the degrees of $g(x)$ and $h(x)$ are ≤ 512, and we can take $d = 2^{10}$. Define the primes $p_1, p_2,$ and p_3 by

$$p_1 = 6946817, \quad p_2 = 7340033, \quad p_3 = 7667713.$$

These primes are between 2^{21} and 2^{23}. Observe that the largest coefficient of $g(x)h(x)$ is no more than

$$512 \cdot 2^{21} \cdot 2^{21} < p_1 p_2 p_3.$$

Hence, using the discrete and fast fourier transforms to compute the coefficients of $f(x)$ modulo each p_j and combining the results with the Chinese Remainder Theorem allows us to determine the coefficients of $f(x) = g(x)h(x)$. Then computing $f(2^{21})$ will give us the value of AB. Note that the same primes (and their corresponding ω's) can be used for every A and B having fewer than 10000 bits.

The LLL Algorithm:

- Background. In 1982, Arjen Lenstra, Hendrik Lenstra, Jr., and László Lovász [9] showed that it is possible to factor a polynomial $f(x) = \sum_{j=0}^{n} a_j x^j \in \mathbb{Z}[x]$ in polynomial time. If n is the degree of $f(x)$ (so $a_n \neq 0$) and H is the height of $f(x)$, that is the maximum of $|a_j|$ for $0 \leq j \leq n$,

then the quantity $n(\log_2 H + \log_2 n + 2)$ can be viewed as an upper bound on the length of the input polynomial $f(x)$. A polynomial time algorithm for factoring $f(x)$ corresponds to an algorithm that runs in time that is polynomial in n and $\log H$. The previous factoring algorithm we described is not polynomial as was seen from the example of Swinnerton-Dyer. The main problem there (which is notably atypical) is that the polynomial $f(x)$ can factor into many small irreducible factors modulo every prime p causing us to have to consider exponentially many possibilities for the mod p reduction of any non-trivial factor of $f(x)$. The algorithm of Lenstra, Lenstra and Lovász (called the LLL-algorithm) is an approach for getting around having to consider all such mod p reductions and thereby provides a polynomial time algorithm for factoring $f(x)$ over the rationals.

- House, Mick Jagger and a general research idea. Mick Jagger had it right, you can't always get what you want. House knew this (see the pilot episode - the song also appeared at the end of the last episode of Season 1 and at the end of the first episode of Season 3). There are many problems in which we want to get the shortest non-zero vector in a lattice (like that of factoring), but we can't or at least no one knows how to. So we can't get what we want. But as Mick Jagger and Cuddy (from House) note, if you try, sometimes you get what you need. In general, if one knows that A leads to a proof of B and A seems unattainable, then it makes sense to try weakening A until it is attainable and seeing if it is enough to imply B (which will likely require a more difficult argument than the one for A implies B but perhaps be based on a similar line of reasoning).

- Lattice background. Let $\mathbf{b}_1, \ldots, \mathbf{b}_n \in \mathbb{Q}^n$, and let $A = (\mathbf{b}_1, \ldots, \mathbf{b}_n)$ be the $n \times n$ matrix with columns $\mathbf{b}_1, \ldots, \mathbf{b}_n$. The lattice $\mathcal{L}$ generated by $\mathbf{b}_1, \ldots, \mathbf{b}_n$ is

$$\mathcal{L} = \mathcal{L}(A) = \mathbf{b}_1 \mathbb{Z} + \cdots + \mathbf{b}_n \mathbb{Z}.$$

We will be interested mainly in the case that $\mathbf{b}_1, \ldots, \mathbf{b}_n$ are linearly independent; in this case, $\mathbf{b}_1, \ldots, \mathbf{b}_n$ is called a basis for $\mathcal{L}$. Observe that given $\mathcal{L}$, the value of $|\det A|$ is the same regardless of the basis $\mathbf{b}_1, \ldots, \mathbf{b}_n$ that is used to describe $\mathcal{L}$. To see this, observe that if $\mathbf{b}'_1, \ldots, \mathbf{b}'_n$ is another basis for $\mathcal{L}$, there are matrices A and B with integer entries such that

$$(\mathbf{b}_1, \ldots, \mathbf{b}_n) A B = (\mathbf{b}'_1, \ldots, \mathbf{b}'_n) B = (\mathbf{b}_1, \ldots, \mathbf{b}_n).$$

Given that $\mathbf{b}_1, \ldots, \mathbf{b}_n$ is a basis for $\mathbb{R}^n$, it follows that AB is the identity matrix and $\det B = \pm 1$. The second equation above then implies

$$|\det (\mathbf{b}'_1, \ldots, \mathbf{b}'_n)| = |\det (\mathbf{b}_1, \ldots, \mathbf{b}_n)|.$$

We set $\det \mathcal{L}$ to be this common value.

- The Gram-Schmidt orthogonaliztion process. Define recursively

$$\mathbf{b}^*_i = \mathbf{b}_i - \sum_{j=1}^{i-1} \mu_{ij} \mathbf{b}^*_j, \qquad \text{for } 1 \leq i \leq n,$$

where

$$\mu_{ij} = \mu_{i,j} = \frac{\mathbf{b}_i \cdot \mathbf{b}^*_j}{\mathbf{b}^*_j \cdot \mathbf{b}^*_j}, \qquad \text{for } 1 \leq j < i \leq n.$$

Then for each $i \in \{1, \ldots, n\}$, the vectors $\mathbf{b}^*_1, \ldots, \mathbf{b}^*_i$ span the same subspace of $\mathbb{R}^n$ (not $\mathcal{L}$) as $\mathbf{b}_1, \ldots, \mathbf{b}_i$. Furthermore, the vectors $\mathbf{b}^*_1, \ldots, \mathbf{b}^*_n$ are linearly independent (hence, non-zero) and pairwise orthogonal (i.e., for distinct i and j, we have $\mathbf{b}^*_i \cdot \mathbf{b}^*_j = 0$). These are easily established.

- Hadamard's inequality. The value of $\det \mathcal{L}$ can be viewed as the volume of the polyhedron with edges parallel to and the same length as $\mathbf{b}_1, \ldots, \mathbf{b}_n$. As indicated by the above remarks, this volume is independent of the basis. Geometrically, it is apparent that

$$\det \mathcal{L} \leq \|\mathbf{b}_1\| \|\mathbf{b}_2\| \cdots \|\mathbf{b}_n\|,$$

where $\|\mathbf{b}_j\|$ is the Euclidean length of $\mathbf{b}_j$ (where "apparent" is limited somewhat to the dimensions we can think in). This is Hadamard's inequality. One can also use the vectors $\mathbf{b}_j^*$ to provide a proof (in any dimensions) as follows. One checks that

$$\det\left(\mathbf{b}_1, \ldots, \mathbf{b}_n\right) = \det\left(\mathbf{b}_1^*, \mathbf{b}_2^* + \mu_{21}\mathbf{b}_1^*, \ldots, \mathbf{b}_n^* + \sum_{j=1}^{n-1}\mu_{nj}\mathbf{b}_j^*\right) = \det\left(\mathbf{b}_1^*, \ldots, \mathbf{b}_n^*\right).$$

Since $\mathbf{b}_1, \ldots, \mathbf{b}_n$ is a basis for $\mathcal{L}$, we deduce that

$$(\det \mathcal{L})^2 = \det\left((\mathbf{b}_1^*, \ldots, \mathbf{b}_n^*)^T(\mathbf{b}_1^*, \ldots, \mathbf{b}_n^*)\right)$$

$$= \det\begin{pmatrix} \|\mathbf{b}_1^*\|^2 & 0 & \cdots & 0 \\ \vdots & & \ddots & \vdots \\ 0 & \cdots & 0 & \|\mathbf{b}_n^*\|^2 \end{pmatrix}$$

$$= \left(\prod_{i=1}^{n} \|\mathbf{b}_i^*\|\right)^2.$$

Thus, $\det \mathcal{L} = \prod_{i=1}^{n} \|\mathbf{b}_i^*\|$. So it suffices to show $\|\mathbf{b}_i^*\| \leq \|\mathbf{b}_i\|$. The orthogonality of the $\mathbf{b}_i^*$'s implies

$$\|\mathbf{b}_i\|^2 = \left\|\mathbf{b}_i^* + \sum_{j=1}^{i-1}\mu_{ij}\mathbf{b}_j^*\right\|^2 = \|\mathbf{b}_i^*\|^2 + \sum_{j=1}^{i-1}\mu_{ij}^2\|\mathbf{b}_j^*\|^2.$$

The sum on the right above is clearly positive so that $\|\mathbf{b}_i^*\| \leq \|\mathbf{b}_i\|$ follows.

- The opposite direction. The Hadamard inequality provides an upper bound on the value of $\det \mathcal{L}$. Hermite proved that there is a constant c_n (depending only on n) such that for some basis $\mathbf{b}_1, \ldots, \mathbf{b}_n$ of $\mathcal{L}$, we have

$$\|\mathbf{b}_1\| \|\mathbf{b}_2\| \cdots \|\mathbf{b}_n\| \leq c_n \det \mathcal{L}.$$

It is known that $c_n \leq n^n$. To clarify a point, Minkowski has shown that there exist n linearly independent vectors $\mathbf{b}_1', \ldots, \mathbf{b}_n'$ in $\mathcal{L}$ such that

$$\|\mathbf{b}_1'\| \|\mathbf{b}_2'\| \cdots \|\mathbf{b}_n'\| \leq n^{n/2} \det \mathcal{L},$$

but $\mathbf{b}_1', \ldots, \mathbf{b}_n'$ is not necessarily a basis for $\mathcal{L}$. Further, we note that the problem of finding a basis $\mathbf{b}_1, \ldots, \mathbf{b}_n$ of $\mathcal{L}$ for which $\|\mathbf{b}_1\| \cdots \|\mathbf{b}_n\|$ is minimal is known to be NP-hard.

- The shortest vector problem. The problem of finding a vector $\mathbf{b} \in \mathcal{L}$ with $\|\mathbf{b}\|$ minimal is not known to be NP-complete, but it may well be. In any case, no one knows a polynomial time algorithm for this problem. (Note that Jeff Lagarias [8] has, however, proved that the problem of finding a vector $\mathbf{b} \in \mathcal{L}$ which minimizes the maximal absolute value of a component is NP-hard.)

Observe that Hermite's result mentioned above implies that there is a constant c_n', depending only on n, such that $\|\mathbf{b}\| \leq c_n' \sqrt[n]{\log \mathcal{L}}$. It is possible for a lattice $\mathcal{L}$ to contain a vector that is much shorter than this, but it is known that the best constant c_n' for all lattices $\mathcal{L}$ satisfies

$$\sqrt{n/(2e\pi)} \leq c_n' \leq \sqrt{n/(e\pi)}.$$

- A lower bound on the shortest vector. The vectors $\mathbf{b}_j^*$ obtained from the Gram-Schmidt orthogonalization process can be used to obtain a lower bound for the shortest vector in a lattice $\mathcal{L}$. More precisely, we have

$$\mathbf{b} \in \mathcal{L}, \ \mathbf{b} \neq \mathbf{0} \implies \|\mathbf{b}\| \geq \min\{\|\mathbf{b}_1^*\|, \|\mathbf{b}_2^*\|, \ldots, \|\mathbf{b}_n^*\|\}. \tag{9}$$

To see this, express $\mathbf{b}$ in the form

$$\mathbf{b} = u_1 \mathbf{b}_1 + \cdots + u_k \mathbf{b}_k, \quad \text{where each } u_j \in \mathbb{Z} \text{ and } u_k \neq 0.$$

Observe that the definition of the $\mathbf{b}_j^*$ imply then that

$$\mathbf{b} = v_1 \mathbf{b}_1^* + \cdots + v_k \mathbf{b}_k^*, \quad \text{for some } v_j \in \mathbb{Q} \text{ with } v_k = u_k.$$

In particular, v_k is a non-zero integer. We deduce that

$$\|\mathbf{b}\|^2 = \left(v_1 \mathbf{b}_1^* + \cdots + v_k \mathbf{b}_k^*\right) \cdot \left(v_1 \mathbf{b}_1^* + \cdots + v_k \mathbf{b}_k^*\right)$$
$$= v_1^2 \|\mathbf{b}_1^*\|^2 + \cdots + v_k^2 \|\mathbf{b}_k^*\|^2 \geq \|\mathbf{b}_k^*\|^2,$$

from which (9) follows.

- Reduced bases. We will want the following important defintion.

Definition: Let $\mathbf{b}_1, \ldots, \mathbf{b}_n$ be a basis for a lattice L and $\mathbf{b}_1^*, \ldots, \mathbf{b}_n^*$ the corresponding basis for $\mathbb{R}^n$ obtained from the Gram-Schmidt orthogonalization process, with μ_{ij} as defined before. Then $\mathbf{b}_1, \ldots, \mathbf{b}_n$ is said to be *reduced* if

(i) $\|\mu_{ij}\| \leq \dfrac{1}{2}$ for $1 \leq j < i \leq n$

(ii) $\|\mathbf{b}_i^* + \mu_{i,i-1} \mathbf{b}_{i-1}^*\|^2 \geq \dfrac{3}{4} \|\mathbf{b}_{i-1}^*\|^2$ for $1 < i \leq n$.

The main work of Lenstra, Lenstra and Lovász [9] establishes an algorithm that runs in polynomial time which constructs a reduced basis of $\mathcal{L}$ from an arbitrary basis $\mathbf{b}_1, \ldots, \mathbf{b}_n$ of $\mathcal{L}$. Our main goal below is to explain how such a reduced basis can be used to factor a polynomial $f(x)$ in polynomial time. We will need to describe the related lattice and an initial basis for it. We begin, however, with some properties of reduced bases.

- A reduced basis contains short vectors. Let $\mathbf{b}_1, \ldots, \mathbf{b}_n$ be a reduced basis for a lattice L and $\mathbf{b}_1^*, \ldots, \mathbf{b}_n^*$ the corresponding basis for $\mathbb{R}^n$ obtained from the Gram-Schmidt orthogonalization process with μ_{ij} as before. The argument for (9) can be modified to show that

$$\mathbf{b} \in \mathcal{L}, \ \mathbf{b} \neq \mathbf{0} \implies \|\mathbf{b}_1\| \leq 2^{(n-1)/2} \|\mathbf{b}\|. \tag{10}$$

In particular, the above inequality holds for the shortest vector $\mathbf{b} \in \mathcal{L}$. To prove (10), observe that (i) and (ii) imply

$$\|\mathbf{b}_i^*\|^2 + \frac{1}{4}\|\mathbf{b}_{i-1}^*\|^2 \geq \|\mathbf{b}_i^*\|^2 + \mu_{i,i-1}^2\|\mathbf{b}_{i-1}^*\|^2 = \|\mathbf{b}_i^* + \mu_{i,i-1}\mathbf{b}_{i-1}^*\|^2 \geq \frac{3}{4}\|\mathbf{b}_{i-1}^*\|^2.$$

Hence, $\|\mathbf{b}_i^*\|^2 \geq (1/2)\|\mathbf{b}_{i-1}^*\|^2$. We deduce that

$$\|\mathbf{b}_i^*\|^2 \geq \frac{1}{2^{i-j}}\|\mathbf{b}_j^*\|^2 \quad \text{for } 1 \leq j < i \leq n. \tag{11}$$

Defining k as in the proof of (9) and following the argument there, we obtain $\|\mathbf{b}\|^2 \geq \|\mathbf{b}_k^*\|^2$. Hence,

$$\|\mathbf{b}\|^2 \geq \|\mathbf{b}_k^*\|^2 \geq \frac{1}{2^{k-1}}\|\mathbf{b}_1^*\|^2 \geq \frac{1}{2^{n-1}}\|\mathbf{b}_1^*\|^2 = \frac{1}{2^{n-1}}\|\mathbf{b}_1\|^2,$$

where the last equation makes use of $\mathbf{b}_1^* = \mathbf{b}_1$. Thus, (10) follows.

Recall that

$$\|\mathbf{b}_i\|^2 = \|\mathbf{b}_i^*\|^2 + \sum_{j=1}^{i-1} \mu_{ij}^2\|\mathbf{b}_j^*\|^2.$$

From (i) and (11), we obtain

$$\|\mathbf{b}_i\|^2 \leq \|\mathbf{b}_i^*\|^2 + \frac{1}{4}\sum_{j=1}^{i-1}\|\mathbf{b}_j^*\|^2 \leq \|\mathbf{b}_i^*\|^2 + \frac{1}{4}\sum_{j=1}^{i-1} 2^{i-j}\|\mathbf{b}_i^*\|^2$$

$$= \left(1 + \frac{1}{4}(2^i - 2)\right)\|\mathbf{b}_i^*\|^2 \leq 2^{i-1}\|\mathbf{b}_i^*\|^2.$$

Using (11) again, we deduce

$$\|\mathbf{b}_j\|^2 \leq 2^{j-1}\|\mathbf{b}_j^*\|^2 \leq 2^{i-1}\|\mathbf{b}_i^*\|^2 \quad \text{for } 1 \leq j \leq i \leq n. \tag{12}$$

We show now the following improvement of (10). Let $\mathbf{x}_1, \mathbf{x}_2, \ldots, \mathbf{x}_t$ be t linearly independent vectors in $\mathcal{L}$. Then

$$\|\mathbf{b}_j\| \leq 2^{(n-1)/2} \max\{\|\mathbf{x}_1\|^2, \|\mathbf{x}_2\|^2, \ldots, \|\mathbf{x}_t\|^2\} \quad \text{for } 1 \leq j \leq t. \tag{13}$$

For each $1 \leq j \leq t$, define a positive integer $m(j)$ and integers u_{ji} by

$$\mathbf{x}_j = \sum_{i=1}^{m(j)} u_{ji}\mathbf{b}_i, \quad u_{jm(j)} \neq 0.$$

By reordering the $\mathbf{x}_j$, we may suppose further that $m(1) \leq m(2) \leq \cdots \leq m(t)$. The linear independence of the $\mathbf{x}_j$ implies that $m(j) \geq j$ for $1 \leq j \leq t$. The proof of (9) implies here that

$$\|\mathbf{x}_j\| \geq \|\mathbf{b}_{m(j)}^*\| \quad \text{for } 1 \leq j \leq t.$$

From (12), we deduce

$$\|\mathbf{b}_j\|^2 \leq 2^{m(j)-1}\|\mathbf{b}_{m(j)}^*\|^2 \leq 2^{m(j)-1}\|\mathbf{x}_j\|^2 \leq 2^{n-1}\|\mathbf{x}_j\|^2 \quad \text{for } 1 \leq j \leq t.$$

The inequality in (13) now follows.

Recall that $\det \mathcal{L} = \prod_{i=1}^{n} \|\mathbf{b}_i^*\|$. We obtain from (12) that

$$\prod_{i=1}^{n} \|\mathbf{b}_i\|^2 \leq \prod_{i=1}^{n} 2^{i-1}\|\mathbf{b}_i^*\|^2 \leq 2^{n(n-1)/2} \prod_{i=1}^{n} \|\mathbf{b}_i^*\|^2 = 2^{n(n-1)/2}(\det \mathcal{L})^2.$$

Thus, from Hadamard's inequality, we obtain

$$2^{-n(n-1)/4}\|\mathbf{b}_1\|\,\|\mathbf{b}_2\| \cdots \|\mathbf{b}_n\| \;\leq\; \det \mathcal{L} \;\leq\; \|\mathbf{b}_1'\|\,\|\mathbf{b}_2'\| \cdots \|\mathbf{b}_n'\|$$

for any basis $\mathbf{b}_1', \ldots, \mathbf{b}_n'$ of $\mathcal{L}$. Recall that finding a basis $\mathbf{b}_1', \ldots, \mathbf{b}_n'$ for which the product on the right is minimal is NP-hard. The above implies that a reduced basis is close to being such a basis. We also note that Hermite's inequality mentioned earlier is a consequence of the above inequality.

- The connection with factoring polynomials. Suppose now that we want to factor a non-zero polynomial $f(x) \in \mathbb{Z}[x]$. Let p be a prime, and consider a monic irreducible factor $h(x)$ of $f(x)$ modulo p^k (obtained say through Berlekamp's algorithm and Hensel lifting). Now, let $h_0(x)$ denote an irreducible factor of $f(x)$ in $\mathbb{Z}[x]$ such that $h_0(x)$ is divisible by $h(x)$ modulo p. Note that $h_0(x)$ being irreducible in $\mathbb{Z}[x]$ implies that the content of $h_0(x)$ (the greatest common divisor of its coefficients) is 1. Our goal here is to show how one can determine $h_0(x)$ without worrying about other factors of $f(x)$ modulo p^k (to avoid the difficulty suggested by Swinnerton-Dyer's example).

We describe a lattice for this approach. Let $\ell = \deg h$. We need only consider the case that $\ell < n$. Fix an integer $m \in \{\ell, \ell+1, \ldots, n-1\}$. We will successively consider such m beginning with ℓ and working our way up until we find $h_0(x)$. In the end, m will correspond to the degree of $h_0(x)$; and if no such $h_0(x)$ is found for $\ell \leq m \leq n-1$, then we can deduce that $f(x)$ is irreducible. We associate with each polynomial

$$w(x) = a_m x^m + \cdots + a_1 x + a_0 \in \mathbb{Z}[x],$$

a vector $\mathbf{b} = (a_0, a_1, \ldots, a_m) \in \mathbb{Z}^{m+1}$. Observe that $\|\mathbf{b}\| = \|w(x)\|$. Let $\mathcal{L}$ be the lattice in $\mathbb{Z}^{m+1}$ spanned by the vectors associated with

$$w_j(x) = \begin{cases} p^k x^{j-1} & \text{for } 1 \leq j \leq \ell \\ h(x)x^{j-\ell-1} & \text{for } \ell+1 \leq j \leq m+1. \end{cases}$$

It is not difficult to see that these vectors form a basis. Furthermore, the polynomials associated with the vectors in $\mathcal{L}$ correspond precisely to the polynomials in $\mathbb{Z}[x]$ of degree $\leq m$ that are divisible by $h(x)$ modulo p^k. In particular, if $m \geq \deg h_0$, the vector, say $\mathbf{b}_0$, associated with $h_0(x)$ is in $\mathcal{L}$. Observe that if k is large enough and $\deg h_0 > \ell$, the coefficients of $h(x)$ are "presumably" large and the value of $\|\mathbf{b}_0\|$ is "seemingly" small. We will show that in fact if k is large enough and $m = \deg h_0$ and $\mathbf{b}_1, \ldots, \mathbf{b}_{m+1}$ is a reduced basis for $\mathcal{L}$, then $\widehat{\mathbf{b}}_1 = \mathbf{b}_0$, where $\widehat{\mathbf{b}}_1$ corresponds to the vector obtained by dividing the components of $\mathbf{b}_1$ by the greatest common divisor of these components (i.e., the polynomial associated with $\widehat{\mathbf{b}}_1$ is the polynomial associated with $\mathbf{b}_1$ with its content removed).

- What does $\mathcal{L}$ have to do with $f(x)$? The lattice $\mathcal{L}$ seemingly has little to do with $f(x)$ as its definition only depends on $h(x)$. Fix $h_0(x)$ as above. We show that $h_0(x)$ is the only irreducible

polynomial in $\mathbb{Z}[x]$ which is associated with a short vector in $\mathcal{L}$ (for k large). For this purpose, suppose $g_0(x)$ is an irreducible polynomial in $\mathbb{Z}[x]$ divisible by $h(x)$ but different from $h_0(x)$ and that R is the resultant of $h_0(x)$ and $g_0(x)$. Note that since $h_0(x)$ and $g_0(x)$ are irreducible in $\mathbb{Z}[x]$, we have $R \neq 0$. The definition of the resultant implies that if R is large, then $\|g_0(x)\|$ must be large (since we are viewing $h_0(x)$ as fixed). So suppose R is not large. There are polynomials $u(x)$ and $v(x)$ in $\mathbb{Z}[x]$ such that

$$h_0(x)u(x) + g_0(x)v(x) = R.$$

We wish to take advantage now of the fact that the left-hand side above is divisible by $h(x)$ modulo p^k, but at the same time we want to keep in mind that unique factorization does not exist modulo p^k. Since $h(x)$ is monic of degree $\ell \geq 1$, the left-hand side is of the form $h(x)w(x)$ modulo p^k where we can now easily deduce that every coefficient of $w(x)$ is divisible by p^k. This implies $p^k|R$. Hence, given k is large, we deduce R is large, giving us the desired conclusion that $\|g_0(x)\|$ is large.

The above argument does more. If $m = \deg h_0(x)$ and $\mathbf{b} \in \mathcal{L}$, then viewing $g_0(x) \in \mathcal{L}$ as the polynomial associated with $\mathbf{b}$, we deduce from the above that either $\|g_0(x)\|$ is large or $R = 0$. In the latter case, since $h_0(x)$ is irreducible and $\deg g_0 \leq m = \deg h_0$, we obtain that $\widehat{\mathbf{b}} = \mathbf{b}_0$.

- How large is large? We take $\mathbf{b} = \mathbf{b}_1$ above (i.e., $g(x)$ is the polynomial associated with the first vector in a reduced basis). Recall that Landau's inequality gives

$$\|h_0(x)\| \leq 2^m \|f(x)\|.$$

On the other hand, by considering the vectors associated with $g_0(x)$ (that is, $\mathbf{b}_1$) and $h_0(x)$ in $\mathcal{L} \subseteq \mathbb{Z}^{m+1}$, we deduce from (10) that

$$\|g_0(x)\| \leq 2^{m/2} \|h_0(x)\|.$$

Thus,

$$\|g_0(x)\| \leq 2^{3m/2} \|f(x)\|.$$

We want this bound on $\|g_0(x)\|$ to assure that $R = 0$ so that $\widehat{\mathbf{b}} = \mathbf{b}_0$.

To see how large p^k needs to be, we recall the Sylvester form of the resultant [17]. We are interested in the resultant R of $g_0(x)$ and $h_0(x)$ where we may suppose that $\deg g_0 \leq m$ and $\deg h_0 \leq m$. From Hadamard's inequality and Landau's inequality, we deduce

$$|R| \leq \|g_0(x)\|^m \|h_0(x)\|^m \leq \|g_0(x)\|^m \left(2^m \|f(x)\| \right)^m = 2^{m^2} \|g_0(x)\|^m \|f(x)\|^m.$$

Our upper bound on $\|g_0(x)\|$ now implies

$$|R| \leq 2^{5m^2/2} \|f(x)\|^{2m}.$$

Hence, we see that if

$$p^k > 2^{5m^2/2} \|f(x)\|^{2m},$$

then the vector $\mathbf{b}_1$ in a reduced basis for the lattice $\mathcal{L}$, where $m = \deg h_0(x)$, is such that the polynomial corresponding to $\widehat{\mathbf{b}}_1$ is $h_0(x)$.

Resources and Further Reading

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