M4: Constructive Mathematics
Lecture 0: What is constructive mathematics?

Patrick E. Farrell

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For some mathematical problems, we can just write down the solution:

$$\text{For } a_0, a_1 \in \mathbb{R}, a_1 \neq 0, \text{ find } x \in \mathbb{R} \text{ such that } a_1 x + a_0 = 0.$$

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For most problems, we can't just write down the solution:

$$\text{For } a_0, \ldots, a_5 \in \mathbb{R}, \text{ find } x \in \mathbb{C} \text{ such that } a_5 x^5 + a_4 x^4 + \cdots + a_0 = 0.$$

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Theorem (Abel, 1824)

There are polynomials of degree 5 and higher that cannot be solved by radicals (addition, subtraction, multiplication, division, and nth root extraction).

Niels Henrik Abel, 1802–1829

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Response A: prove things *about* the solutions.

We could prove that if x is a root of a polynomial with real coefficients, so is $\bar{x}$. Or we could study Vieta's formulae, that (for example) the product of the roots of an n-th degree polynomial is $(-1)^n a_0/a_n$.

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Response B: devise *algorithms* for computing the solutions.

Develop a computational procedure that approximates to arbitrary accuracy the roots of our polynomial: *construct* a sequence that converges to the roots.

The central topic of constructive mathematics is algorithms.

Definition (Algorithm, informal)

An algorithm is a finite set of instructions for solving a mathematical problem. To each input, it associates a sequence of elementary computational steps to calculate some desired output.

The formalisation of this definition is studied in computer science, e.g. with *Turing machines*.

Muḥammad ibn Mūsā al-Khwārizmī, c. 780–850

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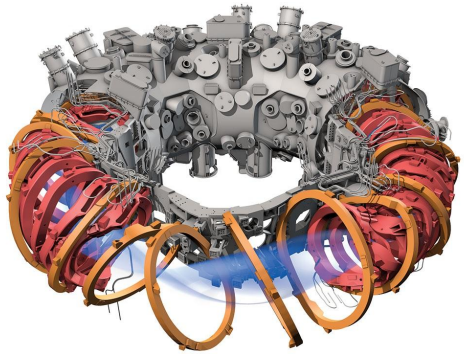
You will see another example in Part A Differential Equations: you will prove that under certain conditions a unique solution exists to the problem

$$\text{find } y(t) \text{ such that } \frac{\mathrm{d}y}{\mathrm{d}t} = f(y,t), \quad y(0) = y_0,$$

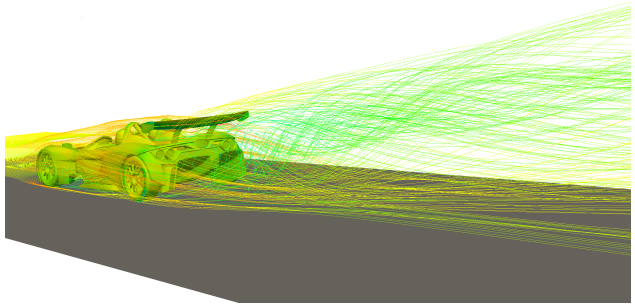
by constructing a sequence of approximations y_n that converges $y_n \to y$.

In applied mathematics, algorithms are used to solve problems arising in science and engineering.

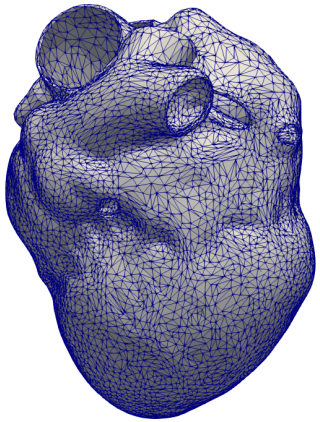
climate



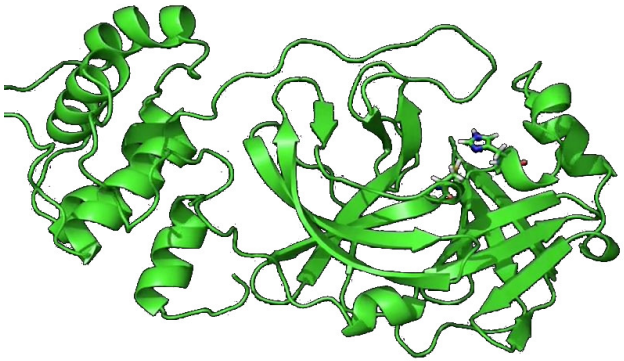
energy



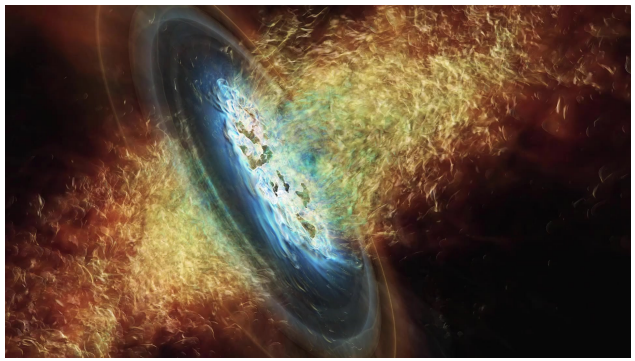
aerodynamics



physiology



covid



galaxies

Questions we ask:

Does our algorithm terminate?

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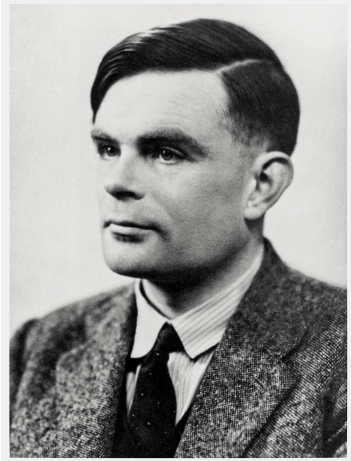
Theorem (Halting problem, 1936)

No algorithm exists that always correctly decides if another algorithm terminates on a given input.



Alan Turing, 1912–1954

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Does our algorithm give the correct answer, and if so, when?

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In later lectures we will see Newton's method for finding a solution x of a general rootfinding problem $f(x) = 0$.

This converges if we start the iteration close to x, but diverges if we start far away.

Isaac Newton, 1643–1727

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How fast does the algorithm converge to the right answer?

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Consider two formulae for π:

$$\pi = 4\sum_{k=0}^{\infty}\frac{(-1)^k}{2k+1}, \qquad \pi^{-1} = \frac{2\sqrt{2}}{99^2}\sum_{k=0}^{\infty}\frac{(4k)!}{k!^4}\frac{26390k+1103}{396^{4k}}.$$

Gottfried Leibniz,
1646–1716

If we approximate the series by its partial sums, how many terms do we require for accuracy to ten digits?

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About 5 billion, vs 2!

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There are many algorithms for sorting a list of n numbers.

The number of comparisons required by a naïve algorithm called *bubble* sort scales like n^2, while the *merge* sort of von Neumann in 1945 scales like $n \log n$. This is much, much faster for large n.

John von Neumann, 1903–1957

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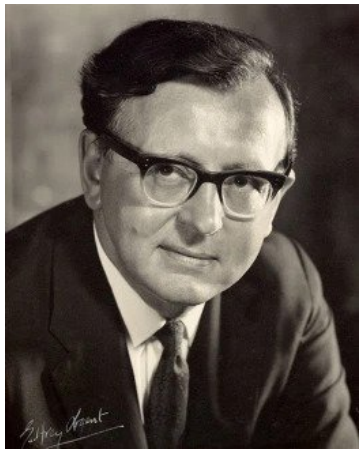
$$p(x) = x^{20} - 210x^{19} + 20615x^{18} + \cdots + 20!.$$



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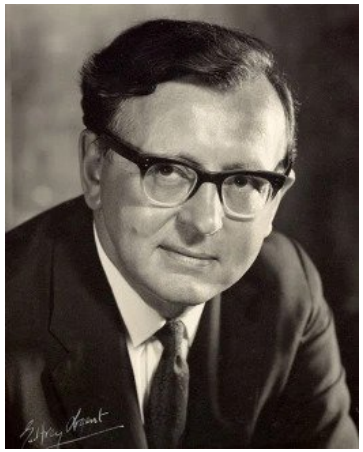

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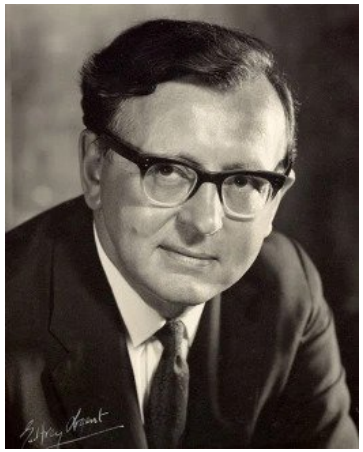

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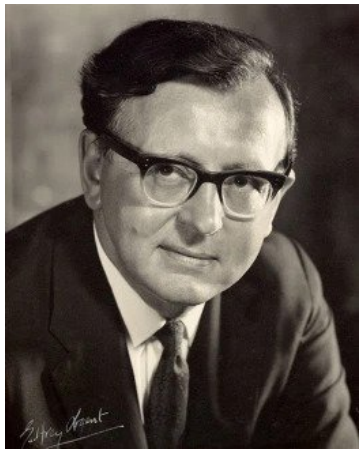

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$10.09527 \pm 0.64350i$	$11.79363 \pm 1.65233i$	$13.99236 \pm 2.51883i$	$16.73074 \pm 2.81262i$	$19.50244 \pm 1.94033i$

M4: Constructive Mathematics
Lecture 1: Euclid's algorithm

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Iterated division

We start with the natural numbers

$$\mathbb{N} = \{0, 1, 2, 3, \dots\},$$

and consider dividing one natural number t by another $b \neq 0$:

$$t = qb + r, \quad r < b.$$

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The game ends when $r = 0$. We're interested in the last remainder before hitting 0. This is the greatest common divisor of the two inputs!

Euclid's method

Here is the *algorithm*. It computes the *greatest common divisor* (also called *highest common factor*) of two numbers.

---

function gcd(t, b)
    $r \leftarrow t \bmod b$
    while $r \neq 0$ **do**
        $t \leftarrow b$
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Note that this algorithm calls another one (the division algorithm).

Theorem (Elements, book VII, c. 300 BCE)

Given any $t, b \in \mathbb{N}$, $0 < b < t$, the Euclidean algorithm computes the greatest common divisor of t and b.

For convenience, let's label each intermediate value:

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\begin{aligned}
t &= q_0 b + r_0 \\
b &= q_1 r_0 + r_1 \\
r_0 &= q_2 r_1 + r_2 \\
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Also for convenience, denote

$$r_{-2} := t, \quad r_{-1} := b.$$

Claim: the algorithm terminates.

Since division yields $r < b$, the sequence of remainders $(r_{-2}, r_{-1}, r_0, \dots)$ is a strictly decreasing sequence of natural numbers. The sequence must therefore eventually reach zero. The algorithm therefore always terminates.

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Let i be the index such that $r_i = 0$.

Claim: r_{i-1} divides r_j, $j < i - 1$ (common divisor).

Since $r_i = 0$, r_{i-1} divides r_{i-2}, i.e.

$$r_{i-2} = q_i r_{i-1}.$$

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Plugging this into the previous iteration tells us that r_{i-1} also divides r_{i-3}:

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Proceeding by induction shows that r_{i-1} divides all remainders in the sequence. In particular, r_{i-1} is a common divisor of the original t and b.

Claim: r_{i-1} is the greatest common divisor.

Assume $d \in \mathbb{N}$ also divides t and b, so there exist $\alpha, \beta \in \mathbb{N}$ such that

$$t = \alpha d, \quad b = \beta d.$$

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The next equation is $b = q_1 r_0 + r_1$, but since d divides both b and r_0, it must also divide r_1. Proceeding by induction, d must divide *all* remainders, including r_{i-1}.

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Thus $d \leq r_{i-1}$, and r_{i-1} is the greatest common divisor of t and b.

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But it is possible to prove a tighter bound!

Theorem

If Euclid's algorithm requires N iterations for a pair of natural numbers $t > b > 0$, the smallest values of t and b for which this is true are the Fibonacci numbers $t = F_{N+2}$ and $b = F_{N+1}$.

Theorem (Complexity of Euclid's algorithm, 1844)

The number of steps taken in Euclid's algorithm can never be more than five times the number of decimal digits of b.

Gabriel Lamé, 1795–1870

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Theorem (Complexity of Euclid's algorithm, 1844)

The number of steps taken in Euclid's algorithm can never be more than five times the number of decimal digits of b.

This result shows that the cost grows *logarithmically* in the size of the input b.

Gabriel Lamé, 1795–1870

Section 2

Diophantine equations

A *Diophantine* equation is an algebraic equation for which solutions are sought in the integers $\mathbb{Z} = \{\ldots, -2, -1, 0, 1, 2, \ldots\}$. They are named after Diophantus of Alexandria (c. 200–290).

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An example of a *linear* Diophantine equation: $4x = 1$. This has no solutions.

Another example: $48x - 35y = 1$. This does indeed have solutions, such as $x = -8, y = -11$. How do we find them?

A *Diophantine* equation is an algebraic equation for which solutions are sought in the integers $\mathbb{Z} = \{\dots, -2, -1, 0, 1, 2, \dots\}$. They are named after Diophantus of Alexandria (c. 200–290).

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Āryabhaṭa (c. 510) realised you could use Euclid's algorithm!

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We make two assumptions about our LDEs:

1. $a > b$. If $a = b$ then we have $a(x + y) = c$, which can either be solved with division or has no solutions. If $a < b$ we could solve $bx' + ay' = c$, and set $(x, y) = (y', x')$.

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2. We assume a and b are coprime—their greatest common divisor is 1. If d divides a and b, then there are two possibilities. If d divides c, we can cancel it from both sides. If d does not divide c, there are no solutions.

Lemma (Bézout's Lemma)

If $\gcd(a, b) = 1$, *then the LDE* $ax + by = 1$ *always has an integer solution.*

Étienne Bézout, 1730–1783

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The statement for integers was already known before Bézout, appearing in the work of Claude Gaspard Bachet de Méziriac in 1624. Bézout's contribution was actually to extend it to polynomials, but his name has stuck to the general principle.

Claude Gaspar Bachet de Méziriac, 1581–1638

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which is the solution $(x, y) = (-8, -11)$ that we saw earlier.

How do we prove Bézout's Lemma? We run Euclid's method *backwards*.

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Proof.

Since $\gcd(a, b) = 1$, we know that iterated divisions of the form

$$\begin{aligned} a &= q_0 b + r_0 \\ b &= q_1 r_0 + r_1 \\ r_0 &= q_2 r_1 + r_2 \\ &\vdots \end{aligned}$$

will eventually reach $r_{i-3} = q_{i-1} r_{i-2} + 1$.

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Let's rewrite this as

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□

We saw in our previous calculations that $48x - 35y = 1$ had a solution $(x, y) = (-8, -11)$. However, there are other solutions, such as $(x, y) = (-43, -59)$. How do we find them *all*? What is the general solution?

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Suppose we have a *particular solution* (x_p, y_p) satisfying $ax_p + by_p = 1$. If we had $(\tilde{x}, \tilde{y})$ such that $a\tilde{x} + b\tilde{y} = 0$, then

$$a(x_p + \tilde{x}) + b(y_p + \tilde{y}) = ax_p + by_p = 1$$

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The general solution to $ax + by = c$ is thus

$$\{c(x_p, y_p) + n(-b, a) : n \in \mathbb{Z}\}.$$

Here is the whole algorithm for solving an LDE $ax + by = c$.

Step 1 Calculate $d = \gcd(a, b)$. If d does not divide c, stop; there are no solutions.

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$$\left\{\hat{c}(x_p, y_p) + n(-\hat{b}, \hat{a}) : n \in \mathbb{Z}\right\}.$$

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$$\begin{aligned} &\{3(-8,-11) + n(35,48) : n \in \mathbb{Z}\} \\ =&\{(-24,-33) + n(35,48) : n \in \mathbb{Z}\}. \end{aligned}$$

Section 3

Euclid for polynomials

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Recall: dividing $p(x)$ by $q(x)$ writes

$$p(x) = c(x)q(x) + r(x)$$

with quotient $c(x)$, and we have either $r(x) = 0$ or $\deg(r) < \deg(q)$.

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A number a is a root of p iff $(x - a)$ divides p, which gives the link between common roots and common divisors.

You've seen polynomial division in school, so we just jump straight into an example of applying Euclid's method. Take

$$p(x) = x^4 + x^3 - 6x^2 + 5x - 1, \quad q(x) = x^3 + x^2 + 3x - 5.$$

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So $(x - 1)$ is the gcd, so $x = 1$ is their only common root:

$$p(1) = 0 = q(1)$$

.

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2. The sequence of remainders yielded by Euclid's method applied to p and p' can be used to compute its *Sturm sequence*. The number of times the Sturm sequence changes sign can be used to calculate how many real roots p has in any given interval (including $(-\infty, \infty)$).

Jacques Charles François Sturm, 1803–1855

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and

$$p_k(x) = \alpha_k(x) \times p_{k-1}(x) + \beta_k \times p_{k-2}(x),$$

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Without specifying α_k or β_k, we can show that p_k and p_{k+1} have no common roots for $k \geq 1$.

Chebyshev polynomials

The main well-conditioned basis for polynomials used in practical computations:

$$T_0(x) = 1, \quad T_1(x) = x,$$
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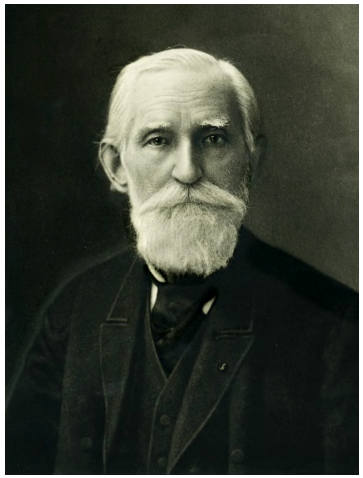

Pafnuty Chebyshev, 1821–1894

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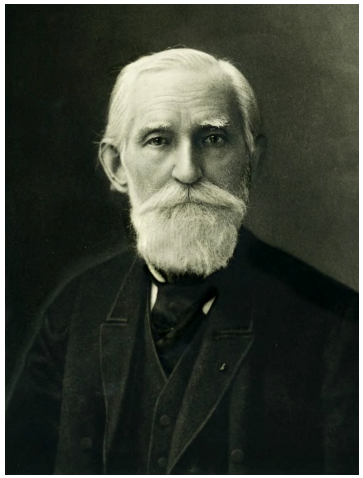

Pafnuty Chebyshev, 1821–1894

Laguerre polynomials

They describe the radial part of the solution of the Schrödinger equation for a one-electron atom:

$$L_0(x) = 1, \quad L_1(x) = -x + 1,$$
$$L_k(x) = \frac{2k+1-x}{k+1} L_{k-1}(x) - \frac{k}{k+1} L_{k-2}(x).$$

Edmond Laguerre, 1834–1886

Proof.

First note that $\deg p_k = k$, by induction.

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This is exactly the division of p_{k+1} by p_k, since $\beta_{k+1} \times p_{k-1}(x)$ is of lower degree than p_k.

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Similarly, p_{k-2} is the remainder on division of p_k by p_{k-1}. Euclid's algorithm thus iterates until it terminates with

$$p_2(x) = \alpha_2(x) \times p_1(x) + \beta_2 p_0(x) = \alpha_2(x) \times x + \beta_2 \times 1,$$

so $\gcd(p_k, p_{k+1})$ is a nonzero constant (no roots). □

M4: Constructive Mathematics
Lecture 2: Rootfinding and fixed points

Patrick E. Farrell

University of Oxford

In the previous lecture we saw that we could use Euclid's method to compute the common roots of two polynomials p and q.

This, however, is very limited. We will want to find roots of general (not necessarily polynomial) functions $f : \mathbb{R} \to \mathbb{R}$.

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For this, we turn to *rootfinding* algorithms. There are many different ones, differing in efficiency, robustness, and applicability.

Rootfinding problem

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Another use: if you want to calculate the decimal expansion of a number (like $\sqrt{2}$), set up a suitable equation, like

$$x^2 - 2 = 0$$

and apply a rootfinding algorithm.

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- Does the algorithm give the correct answer?
- How fast does the algorithm converge to the answer?
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Different algorithms will trade off termination, convergence speed, and operation count.

Section 2

Bisection

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If $f : [a, b] \to \mathbb{R}$ *is continuous with* $f(a)f(b) < 0$*, then there exists* $x^\star \in (a, b)$ *with* $f(x^\star) = 0$*.*

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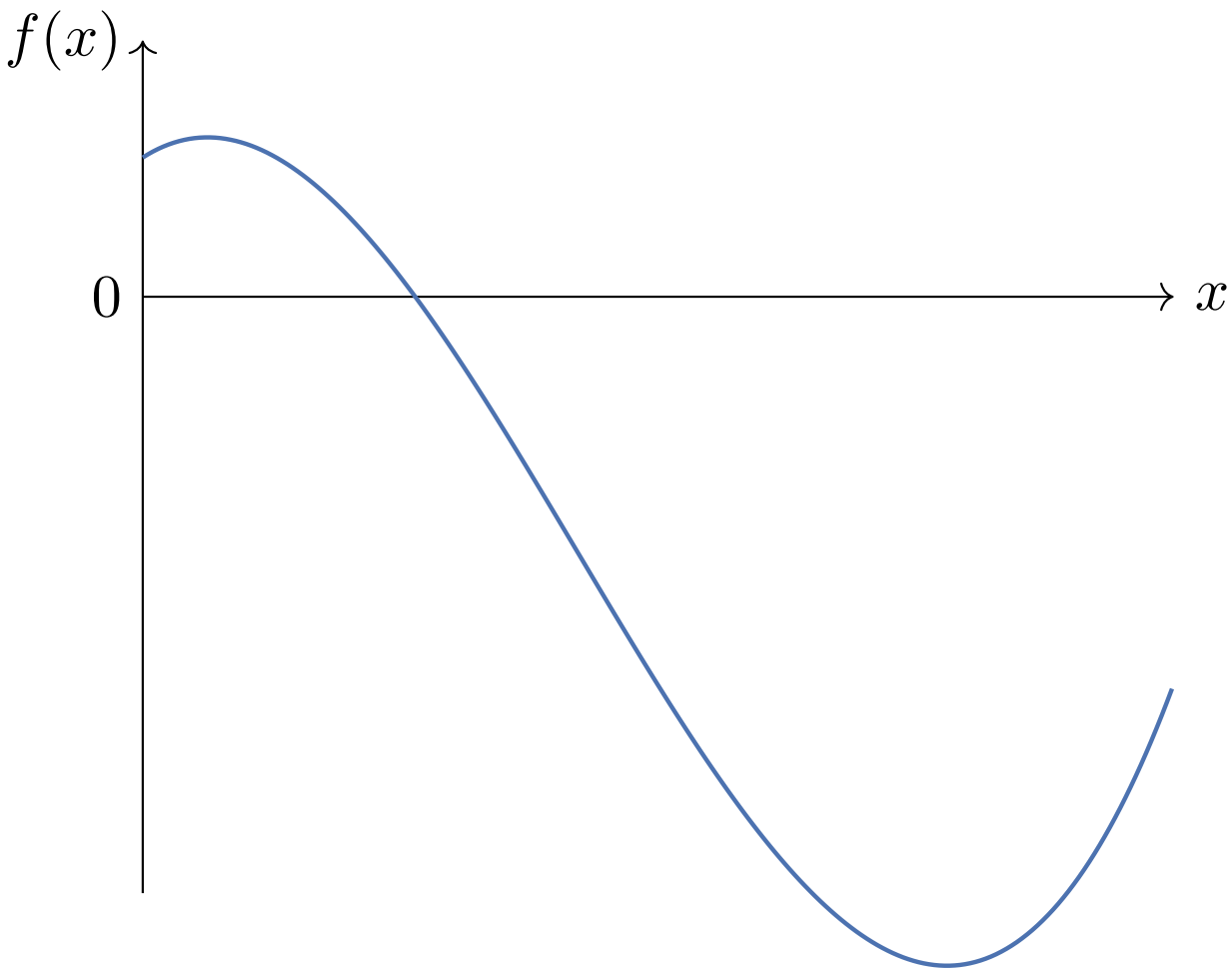
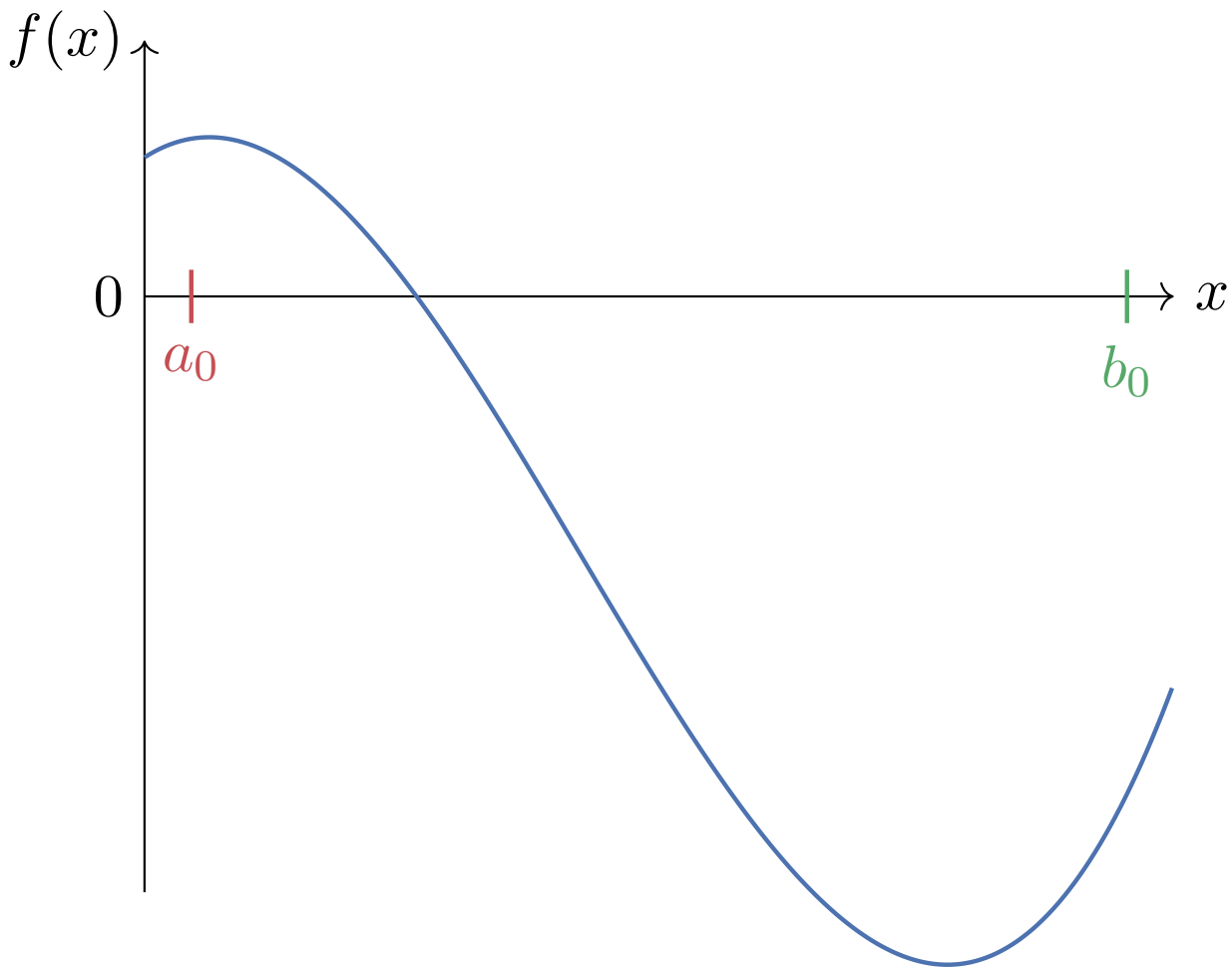
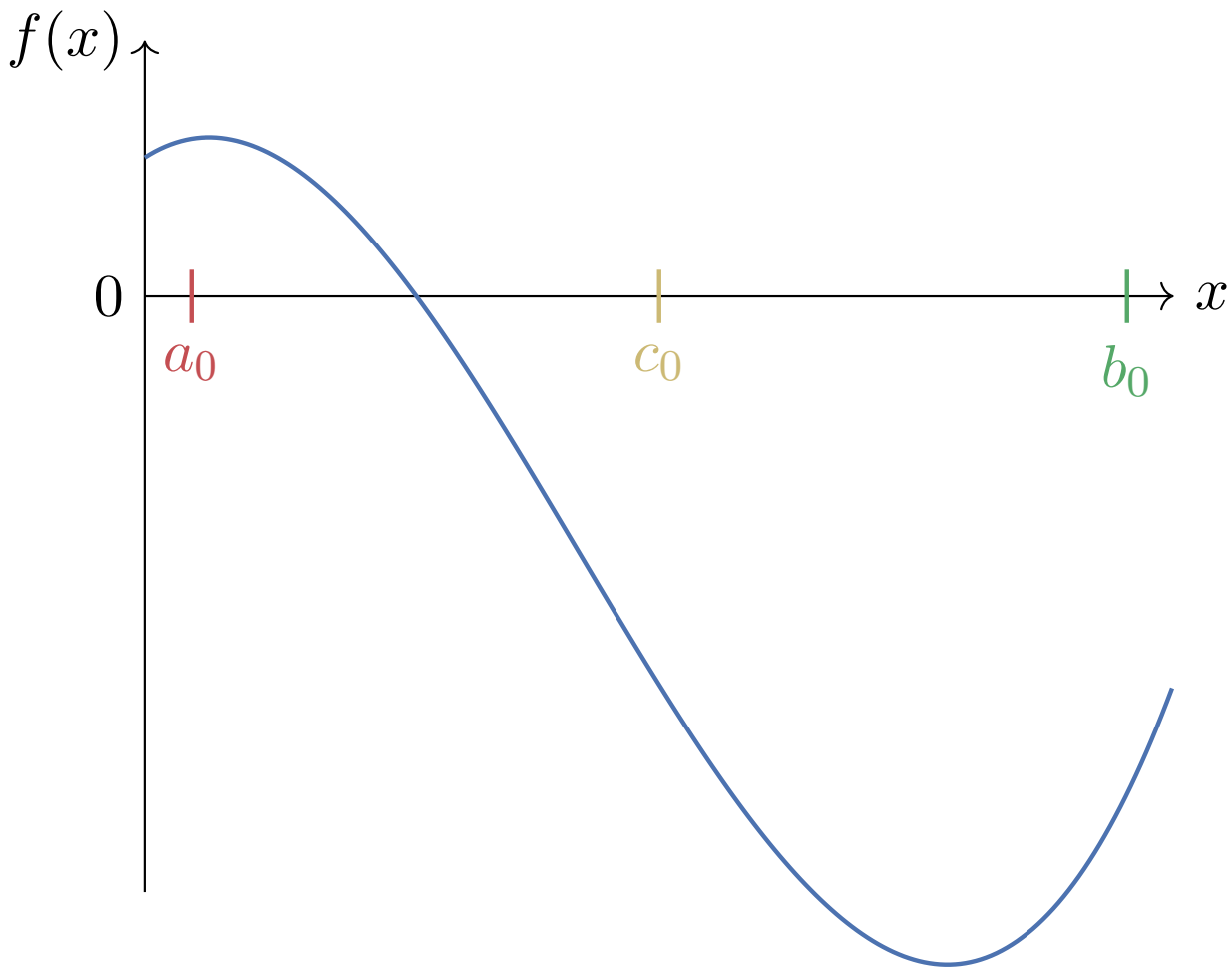
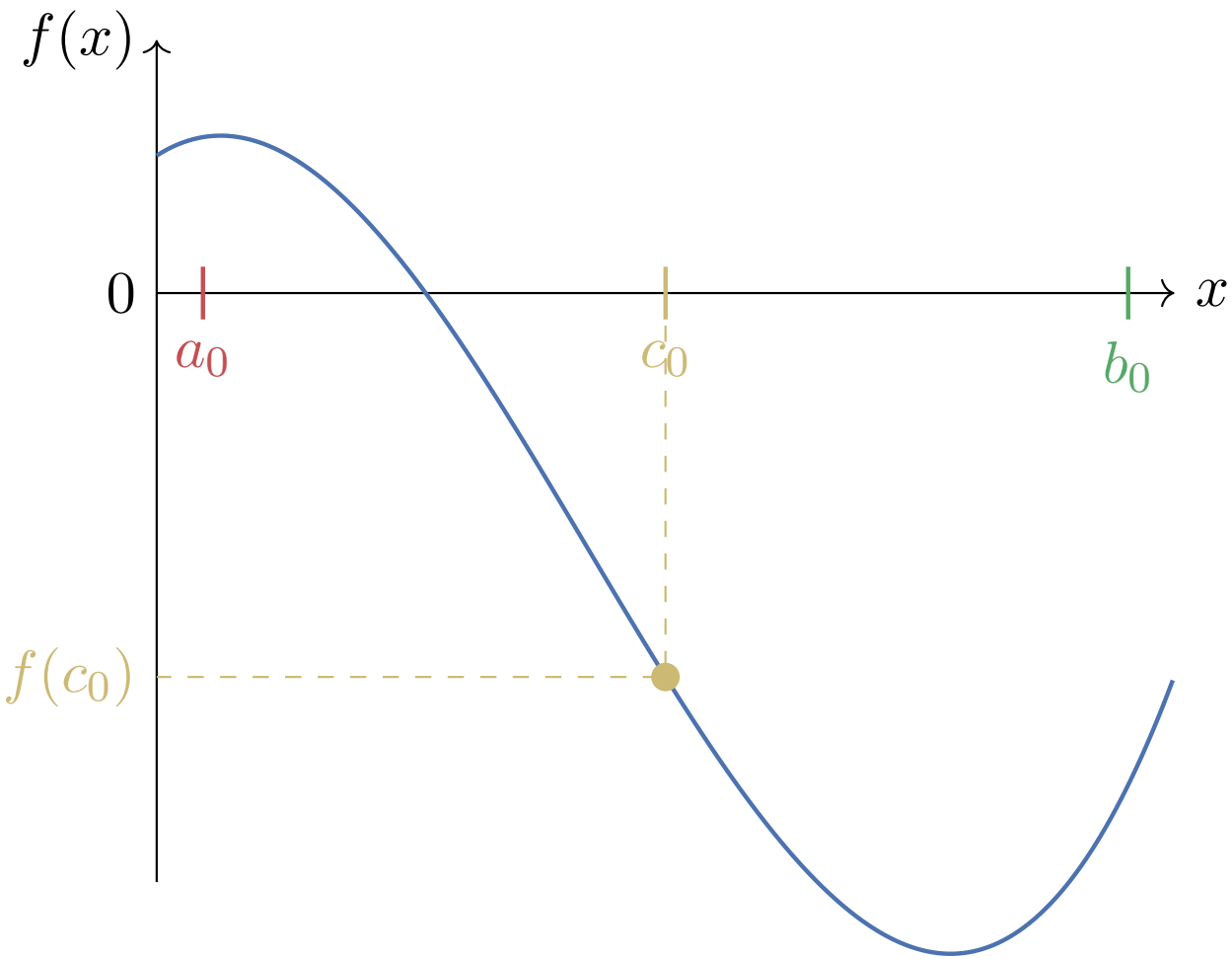
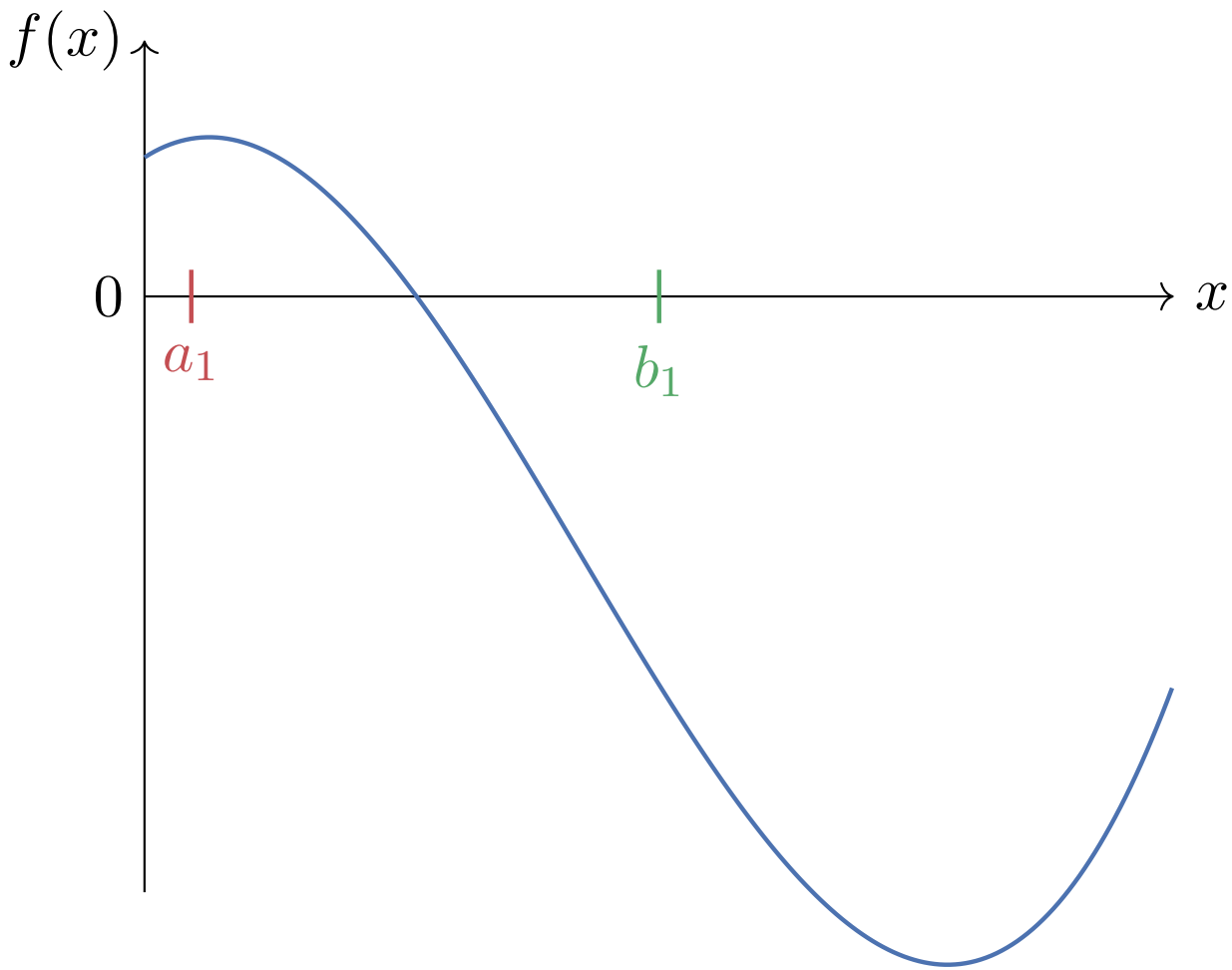
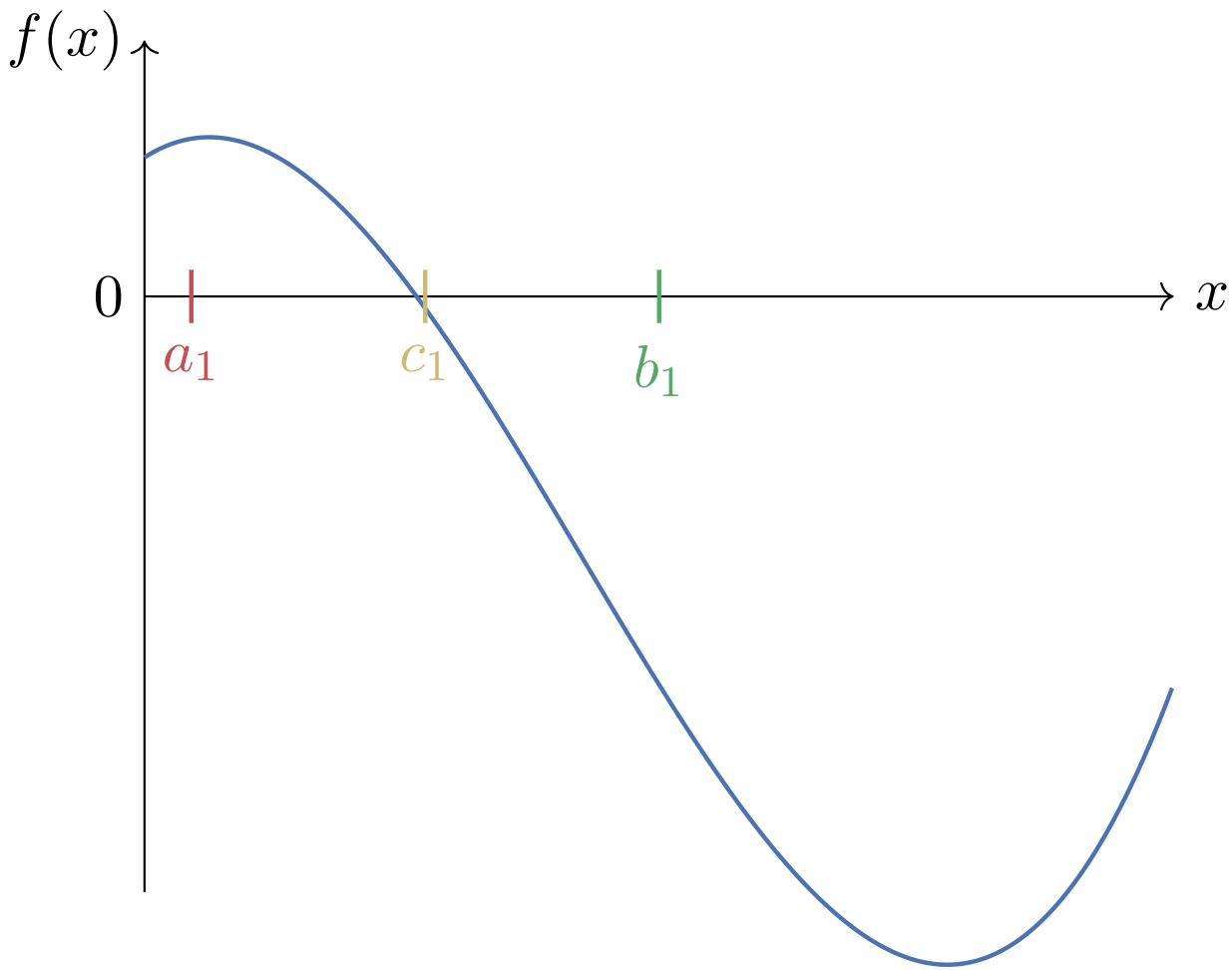

$f(x)$
0
x

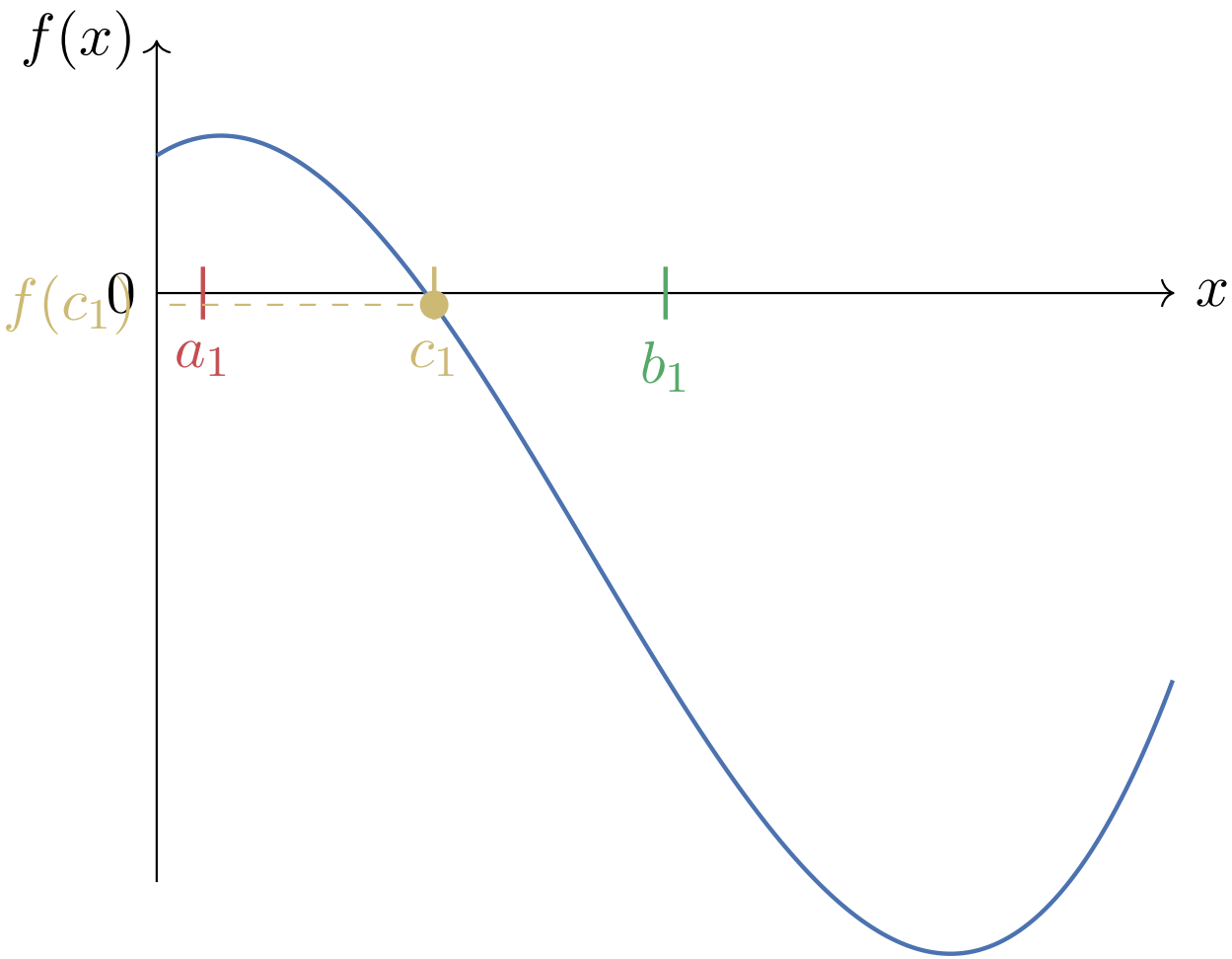

$f(x)$
0
a_0
b_0
x

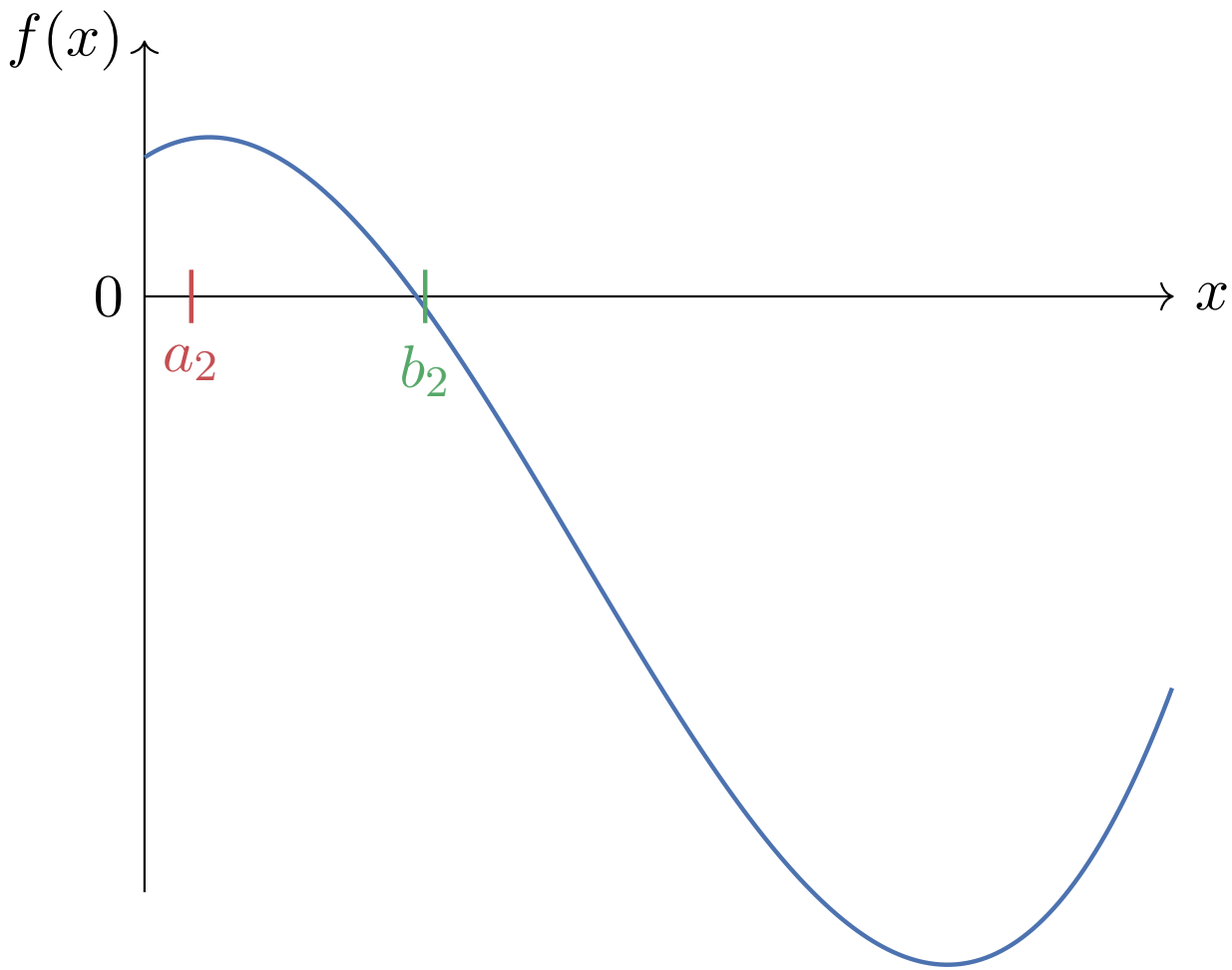

$f(x)$
0
x
a_0
c_0
b_0

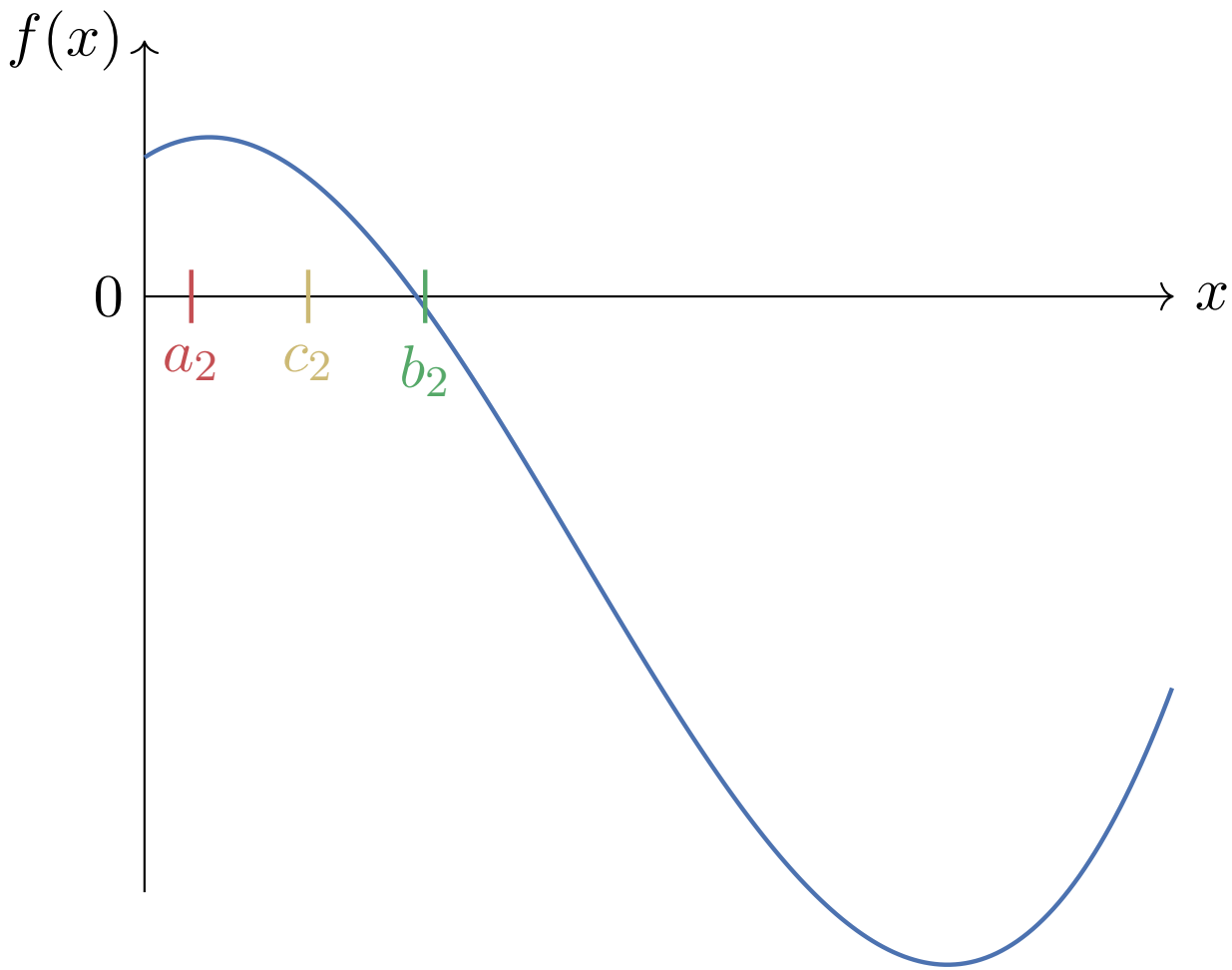

$f(x)$
0
x
a_0
c_0
b_0
$f(c_0)$

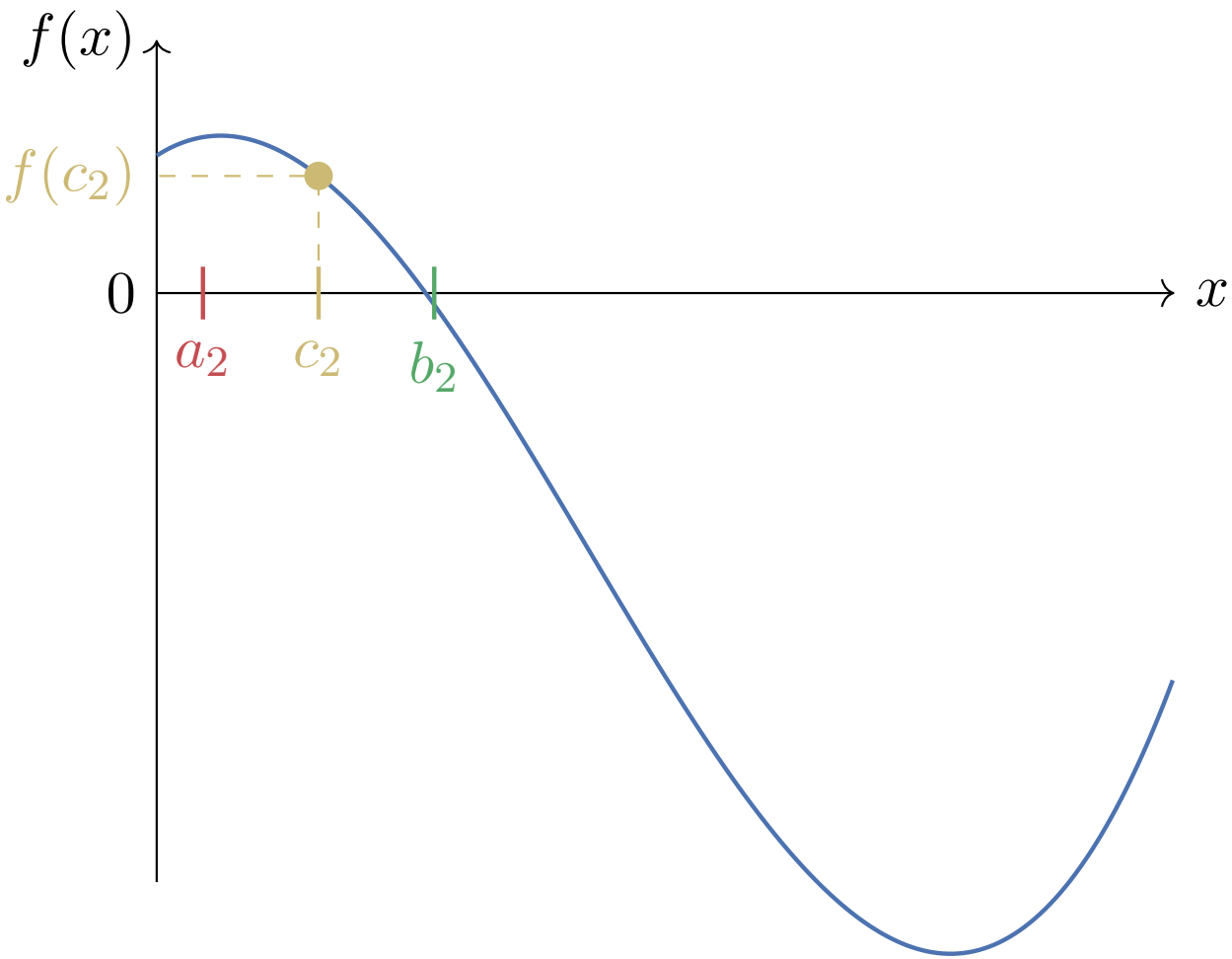

$f(x)$
0
x
a_1
b_1

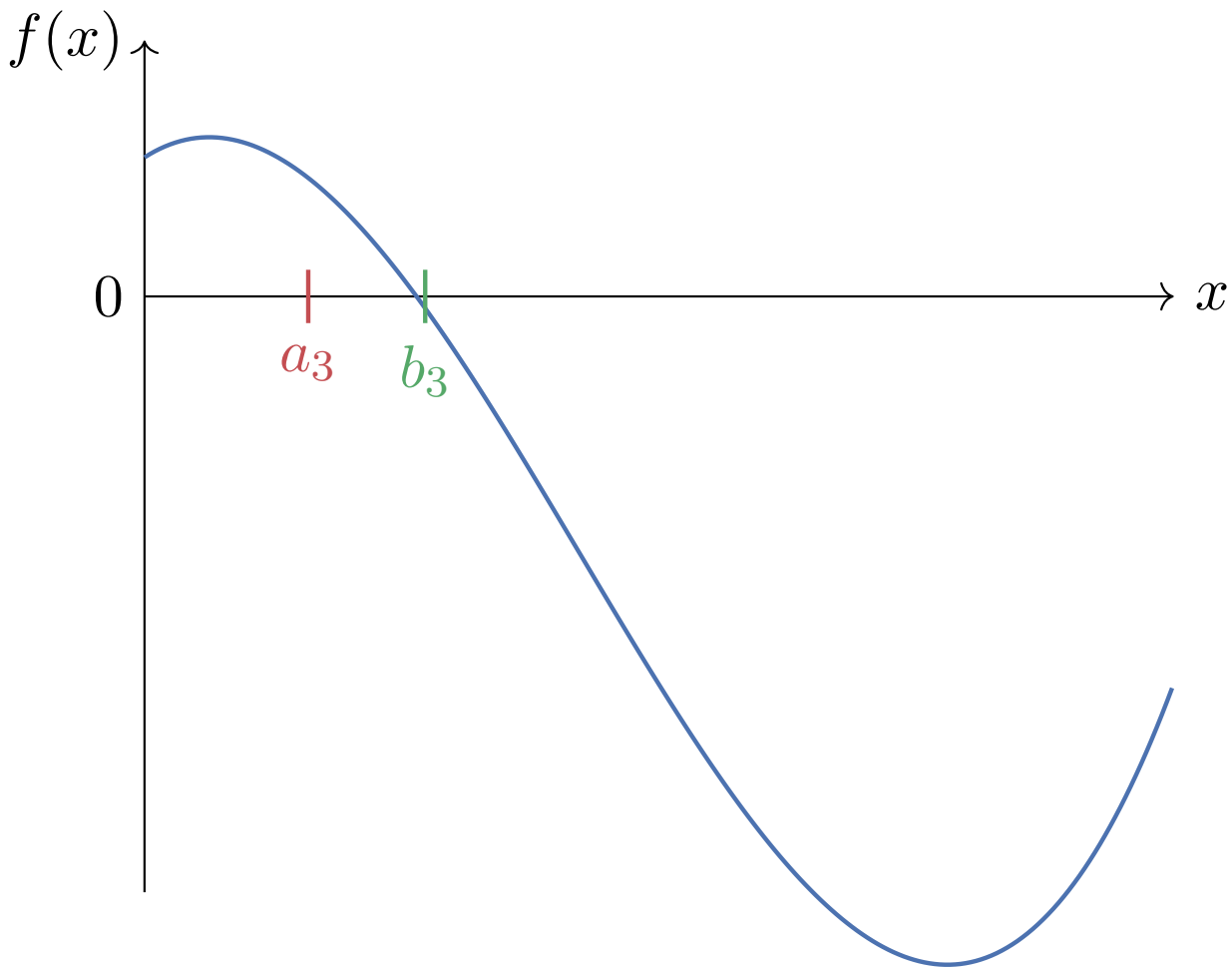

$f(x)$
0
x
a_1
c_1
b_1

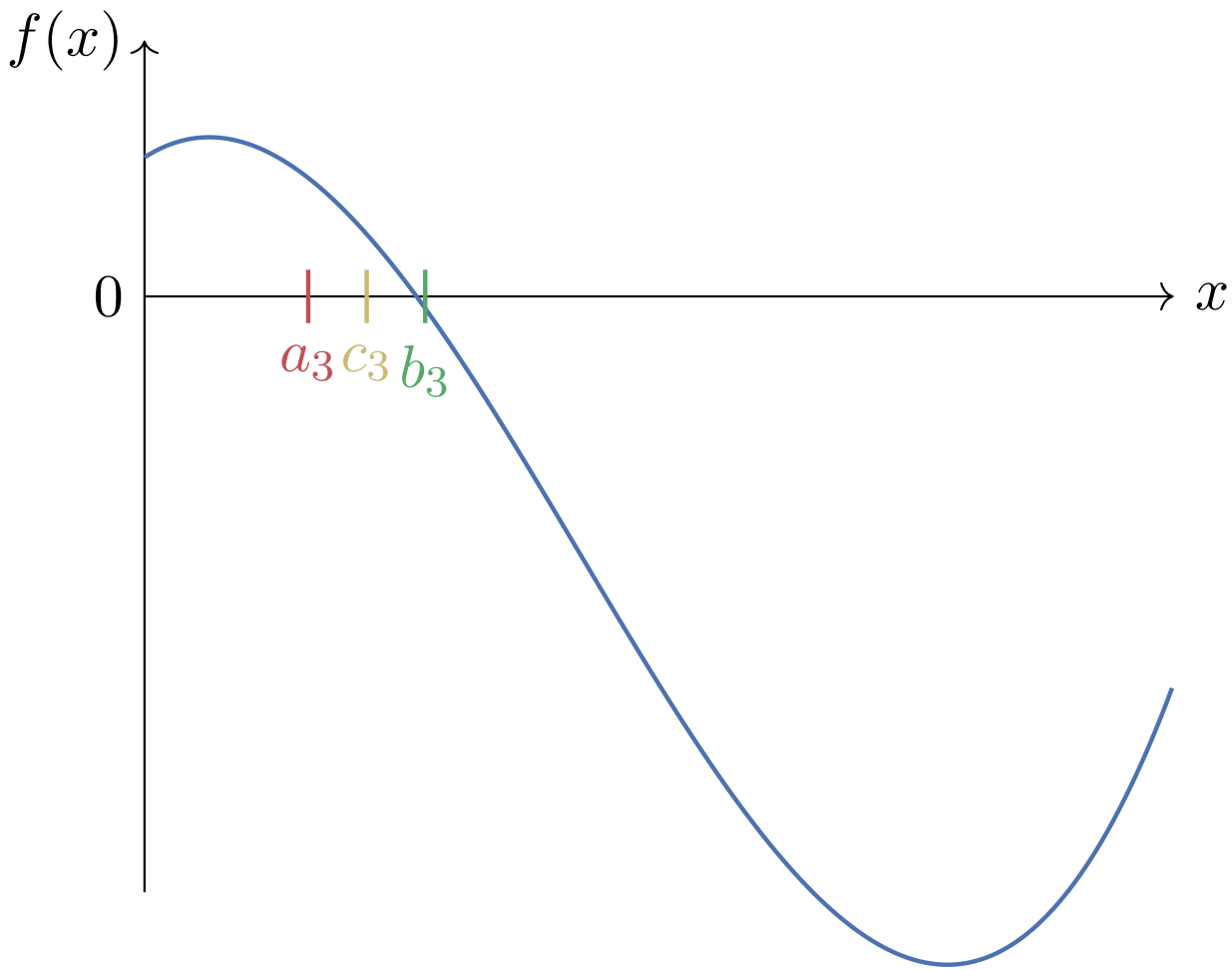

$f(x)$
$f(c_1)$
0
x
a_1
c_1
b_1

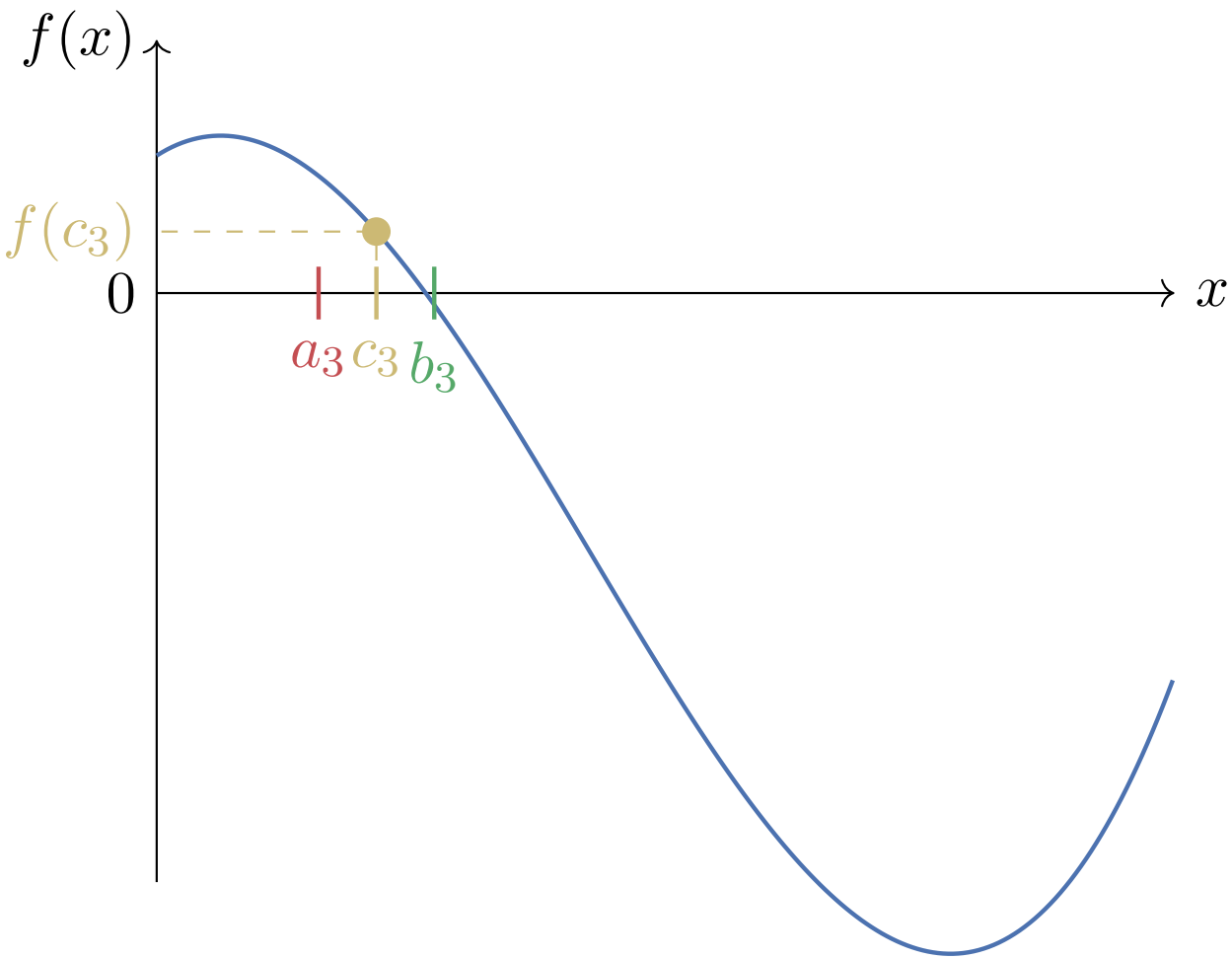

$f(x)$
0
x
a_2
b_2

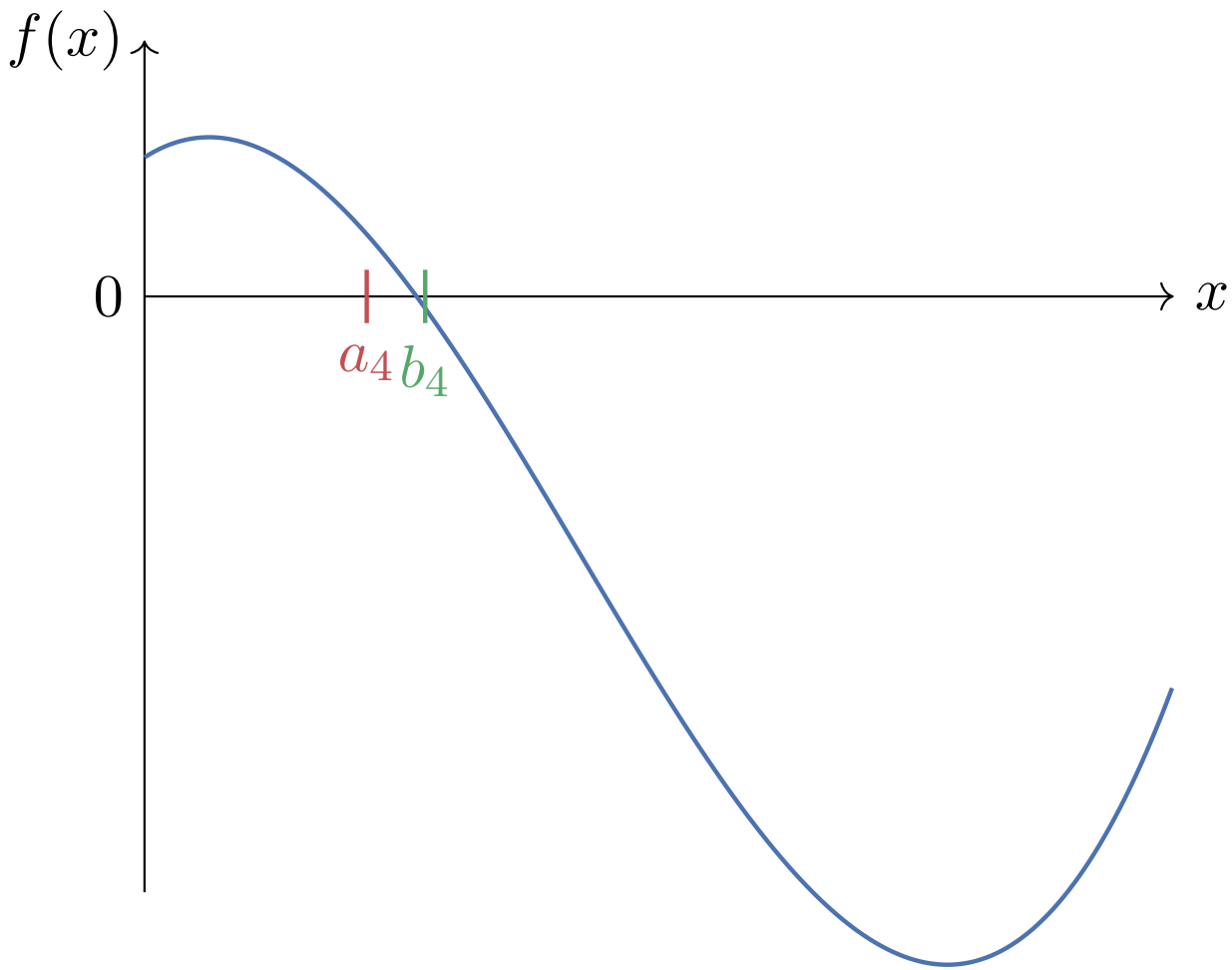

$f(x)$
0
x
a_2
c_2
b_2


$f(x)$
$f(c_2)$
0
a_2
c_2
b_2
x


$f(x)$
0
x
a_3
b_3


$f(x)$
0
x
a_3
c_3
b_3


$f(x)$
$f(c_3)$
0
x
a_3 c_3 b_3


$f(x)$
0
x
a_4
b_4

Let's state this as an algorithm.

Assume $f : [a, b] \to \mathbb{R}$ is continuous, $f(a)f(b) < 0$, and $\mathrm{tol} > 0$.

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Note this only uses the *sign* of the output of $f(x)$.

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Lemma

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Proof.

In the k-th iteration of the while loop, either the function returns or it shrinks $|b-a|$ by a factor of 2. For any $\mathrm{tol} > 0$, there exists $k \in \mathbb{N}$ such that $\mathrm{tol} < |b-a|/2^{k+1}$, so the algorithm must terminate. □

Let's do an example. Let's try to solve $x = \cos x$, so $f(x) = x - \cos x$.

Let's start with $[a, b] = [-10, 10]$. $f(-10) \approx -9.16$, $f(10) \approx 10.83$, so we're good to go.

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The true solution is approximately $x \approx 0.739085$, so we're getting there, slowly.

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- ✗ It can never find roots of even multiplicity.

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- ✗ It can never find roots of even multiplicity.

Definition (Multiplicity of a root)

A root $x^\star$ of $f(x)$ has multiplicity k if $f^{(n)}(x^\star) = 0$ for all $n < k$, and $f^{(k)}(x^\star) \neq 0$.

Comments on bisection

- ✓ When it applies, it is guaranteed to converge.
- ✓ The method is very simple and very robust.
- ✓ Its smoothness requirements are low (only continuity).
- ✗ The interval of interest only reduces by a factor of two each time.
- ✗ It can be hard to find the initial points $[a, b]$ that bracket a root.
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Later we will study other methods with different sets of advantages and disadvantages.

Section 3

Rate of convergence of a sequence

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Suppose $(x_i) \to x^\star$. We say the sequence converges linearly if there exists $\mu \in (0, 1)$ such that

$$\lim_{i \to \infty} \frac{|x_{i+1} - x^\star|}{|x_i - x^\star|} = \mu.$$

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For bisection, the sequence of the midpoints of the intervals converges linearly with $\mu = 1/2$.

Can you go faster?

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Definition (Superlinear convergence of a sequence)

Suppose $(x_i) \to x^\star$. We say the sequence converges superlinearly if

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For example, the sequence

$$(\frac{1}{2^{2^n}}) = (\frac{1}{2}, \frac{1}{4}, \frac{1}{16}, \frac{1}{256}, \frac{1}{65535}, \dots) \to 0$$

has the ratio of successive terms going to zero too.

We can further classify superlinear convergence:

Definition (Order of convergence of a sequence)

Suppose $(x_i) \to x^\star$, superlinearly. The sequence converges with order q if

$$\lim_{i \to \infty} \frac{|x_{i+1} - x^\star|}{|x_i - x^\star|^q} = M$$

for some $M > 0$ (not necessarily $M < 1$).

We call $q = 2$ quadratic convergence, $q = 3$ cubic convergence, etc.

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We will see rootfinding methods with orders of convergence $q = 2$ and $q = 3$. To develop these, we must first understand *fixed point iterations*.

Section 4

Fixed point iterations

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We can translate between rootfinding problems and fixed point problems. For example, if you want to find the fixed points $g(x) = x$, then you can find the roots of $f(x) := g(x) - x$.

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Vice versa, if you have a rootfinding problem $f(x) = 0$, you could search for fixed points of $g(x) := f(x) + x$. There are other ways of transforming between them, of course.

Transforming between the two problems is useful because there are powerful theorems that apply to finding fixed points. There's even a whole course, C4.6 Fixed Point Methods for Nonlinear PDEs, on this subject.

When can we show fixed points exist?

Theorem (Brouwer's fixed point theorem)

If $g : [a, b] \to [a, b]$ is continuous, then it has a fixed point.

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Warning (endomorphism)

Note that g must send $[a, b]$ to $[a, b]$, i.e. is an *endomorphism*. This result does *not* hold for general $g : [a, b] \to \mathbb{R}$, such as $g(x) = x + 1$.

Proof.

Since $g(x) \in [a, b]$, we have $a \leq g(x) \leq b$ for all $x \in [a, b]$. Thus $f(x) := g(x) - x$ has $f(a) \geq 0$ and $f(b) \leq 0$.

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If either inequality is an equality, we have a fixed point. So assume that $f(a) > 0$ and $f(b) < 0$.

A root $x^\star$ of $f(x)$ thus exists in (a, b) by Bolzano's Theorem, with $g(x^\star) = x^\star$. □

That's not all! You can get uniqueness of the fixed point under stronger conditions.

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Theorem

*If $g : [a, b] \to [a, b]$ is differentiable with $|g'(x)| < 1$ for every $x \in (a, b)$, then g has a **unique** fixed point in (a, b).*

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Theorem (Mean value theorem, 1823)

If $g : [a, b] \to \mathbb{R}$ is differentiable, then there exists some $c \in (a, b)$ such that

$$g'(c) = \frac{g(b) - g(a)}{b - a}.$$

Augustin-Louis Cauchy FRS
1789–1857

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Assume without loss of generality that $p < q$. Applying the MVT in $[p, q] \subset [a, b]$, we find that there exists $r \in (p, q)$ such that

$$g'(r) = \frac{g(q) - g(p)}{q - p} = \frac{q - p}{q - p} = 1.$$

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How do we turn this into an algorithm?

Take $x_0 \in [a, b]$ and set $x_{i+1} = g(x_i)$!

Assume $g : [a, b] \to [a, b]$, and $x_0 \in [a, b]$.

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function fixedpoint(g, x_0, tol)
    $x \leftarrow x_0$
    while $|g(x) - x| > \text{tol}$ **do**
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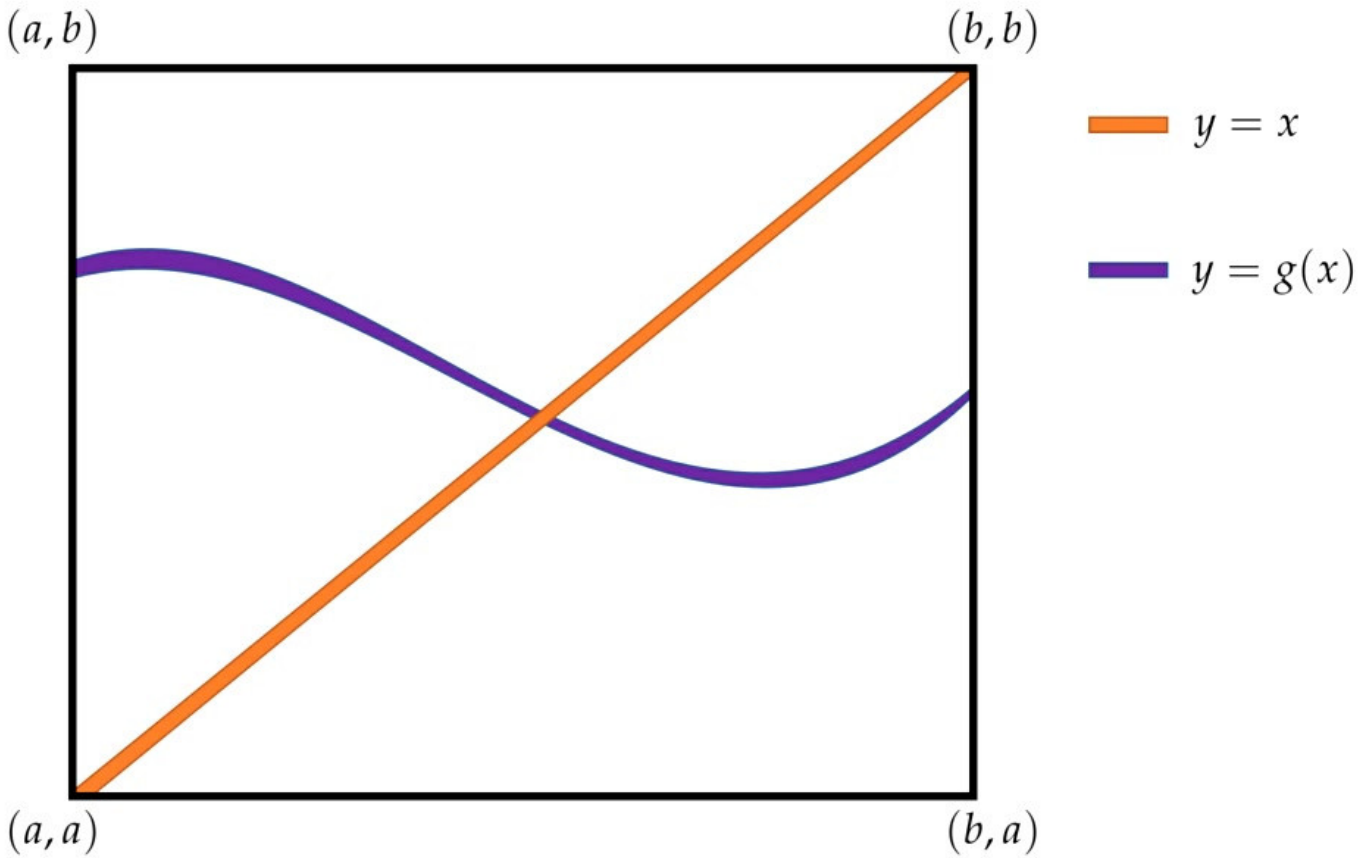
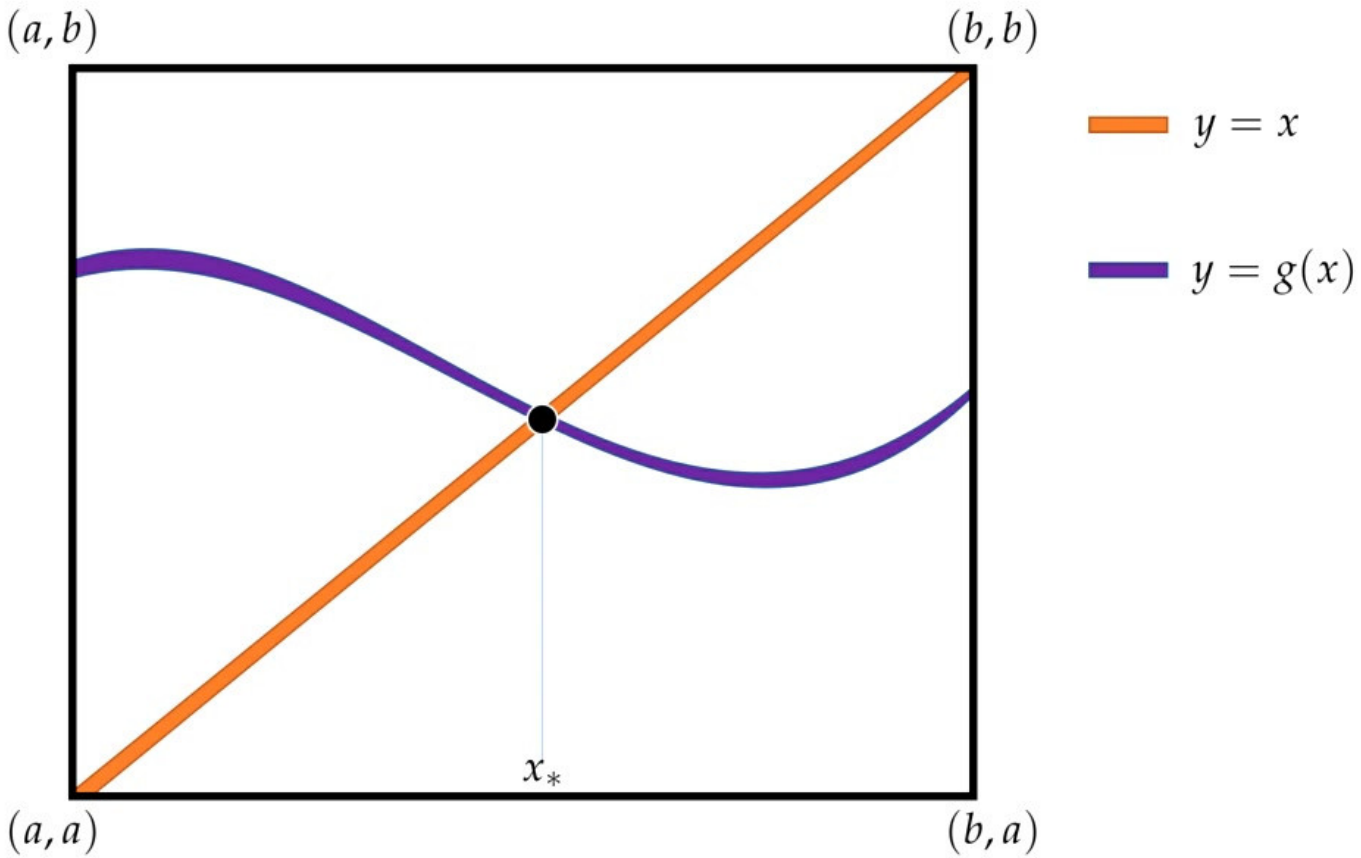
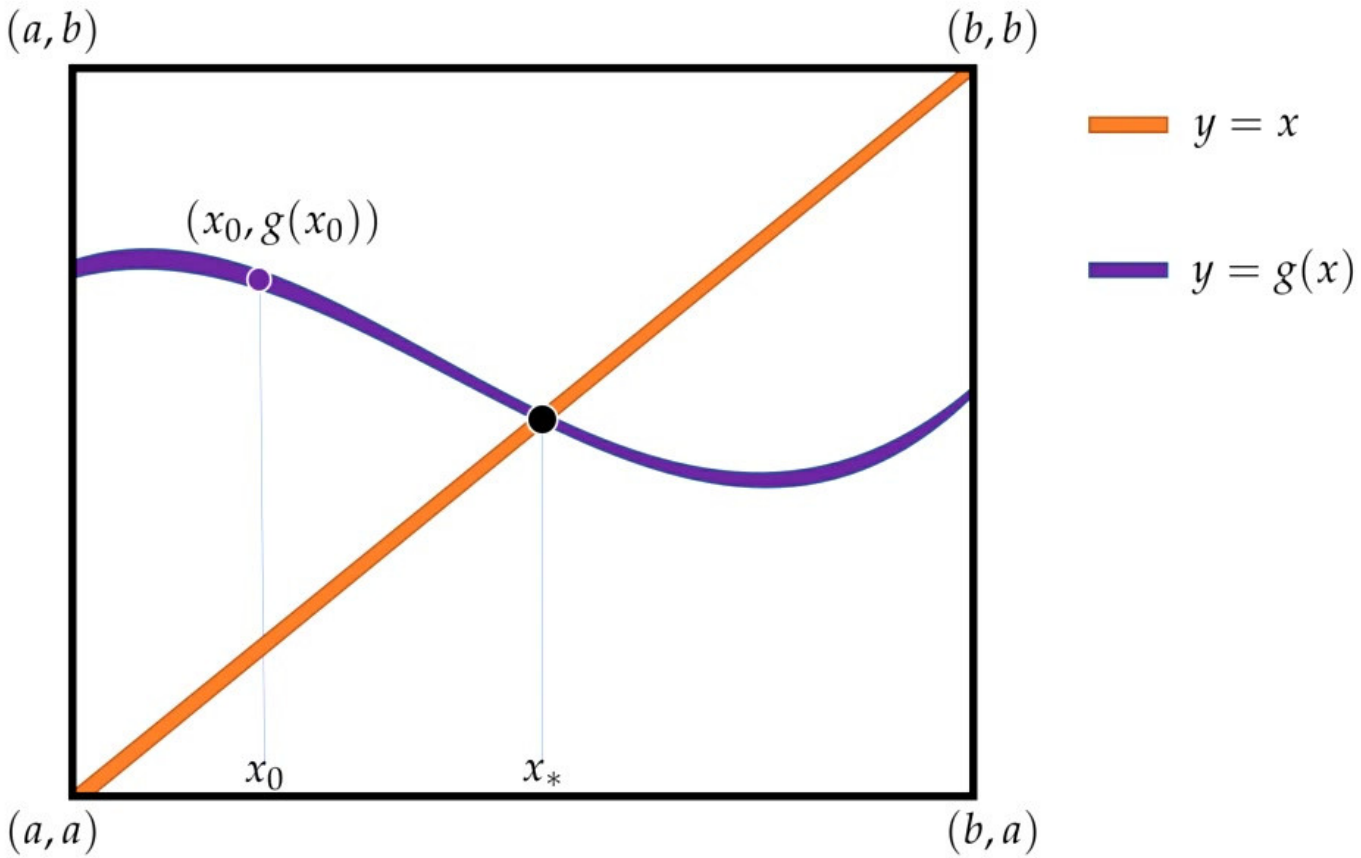
Our goal is to investigate when this converges.


(a, b)
(b, b)
$y = x$
$y = g(x)$
(a, a)
(b, a)


(a, b)
(b, b)
$y = x$
$y = g(x)$
x_*
(a, a)
(b, a)


(a, b)
(b, b)
$y = x$
$(x_0, g(x_0))$
$y = g(x)$
x_0
x_*
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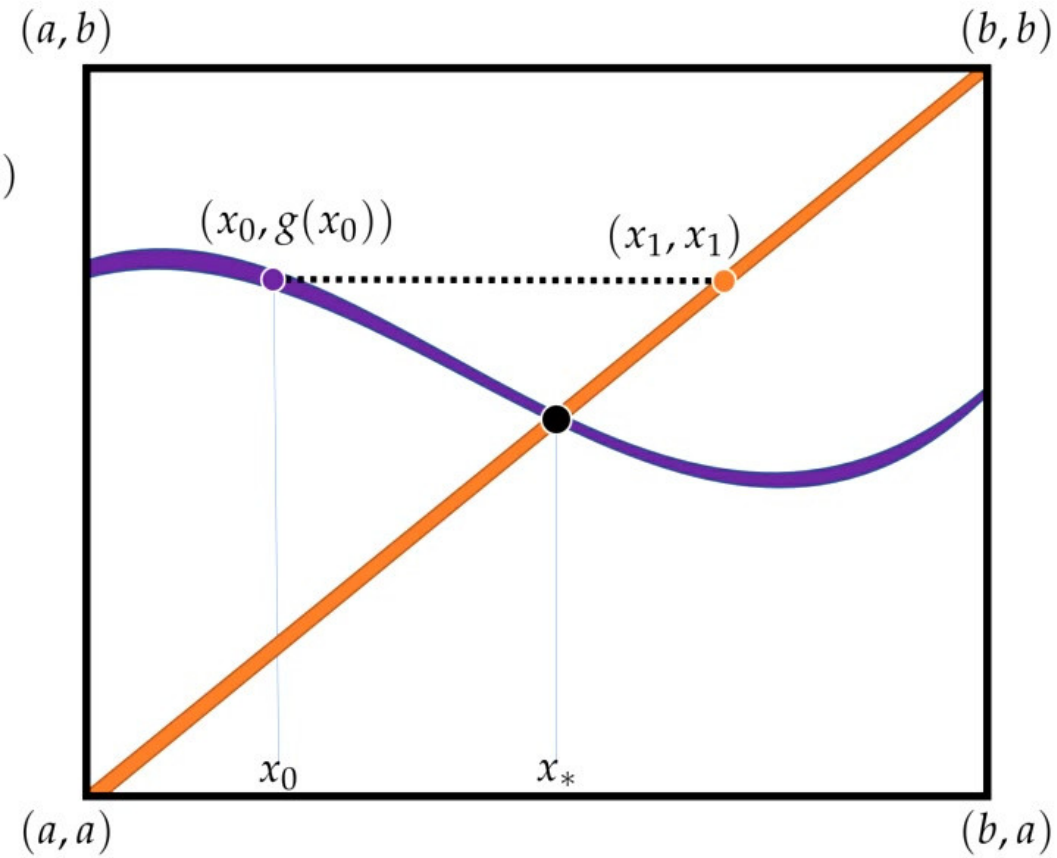

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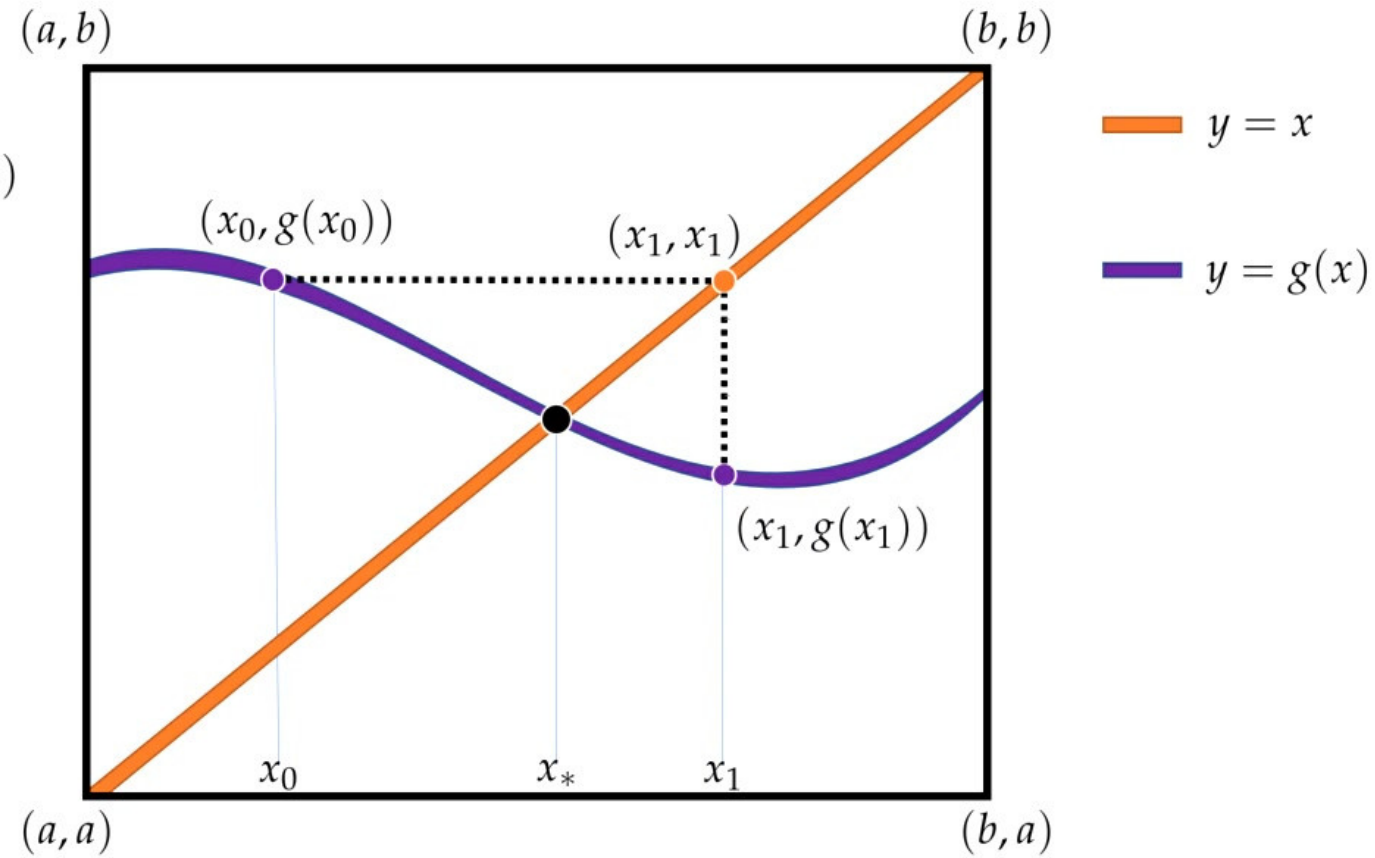

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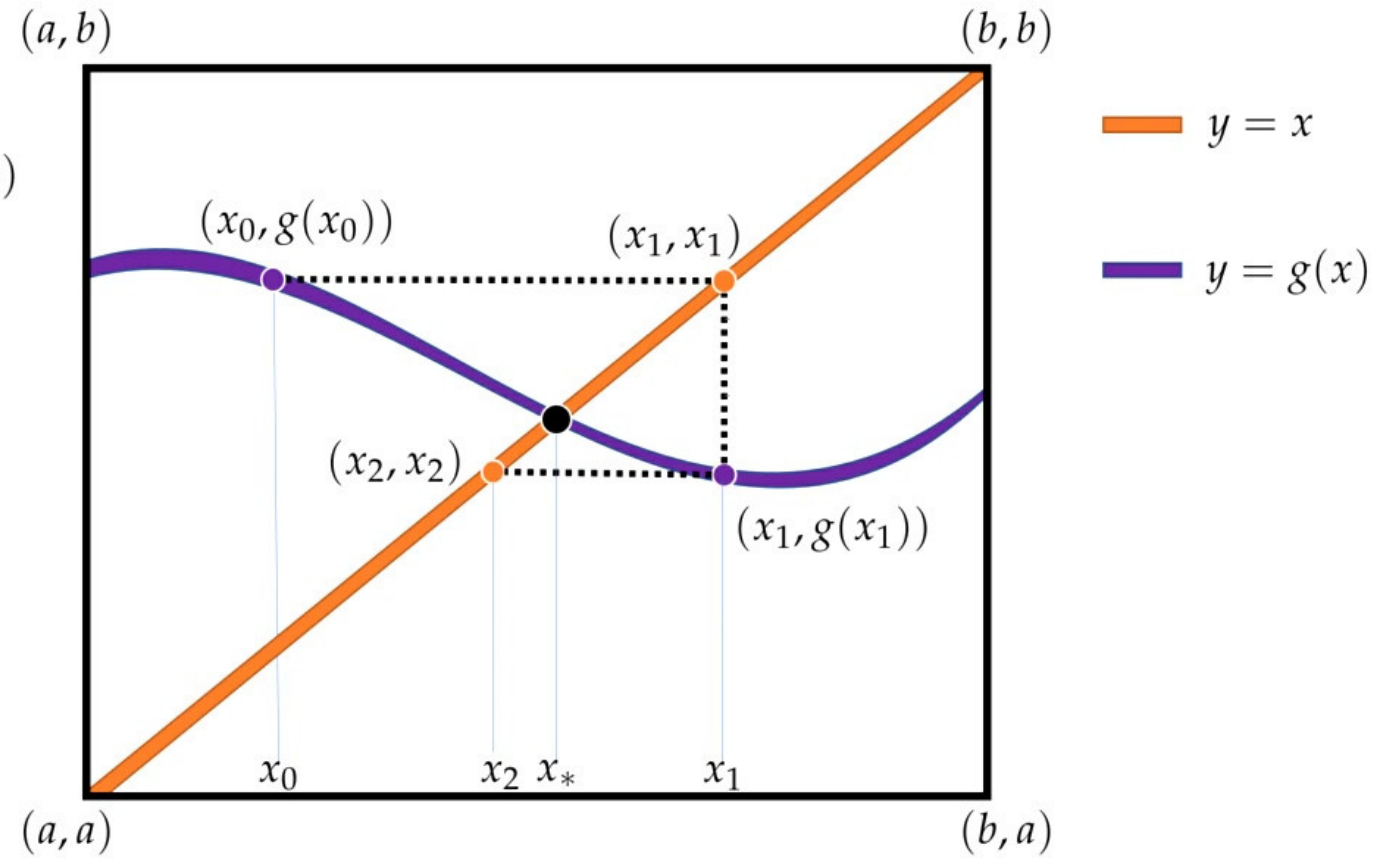

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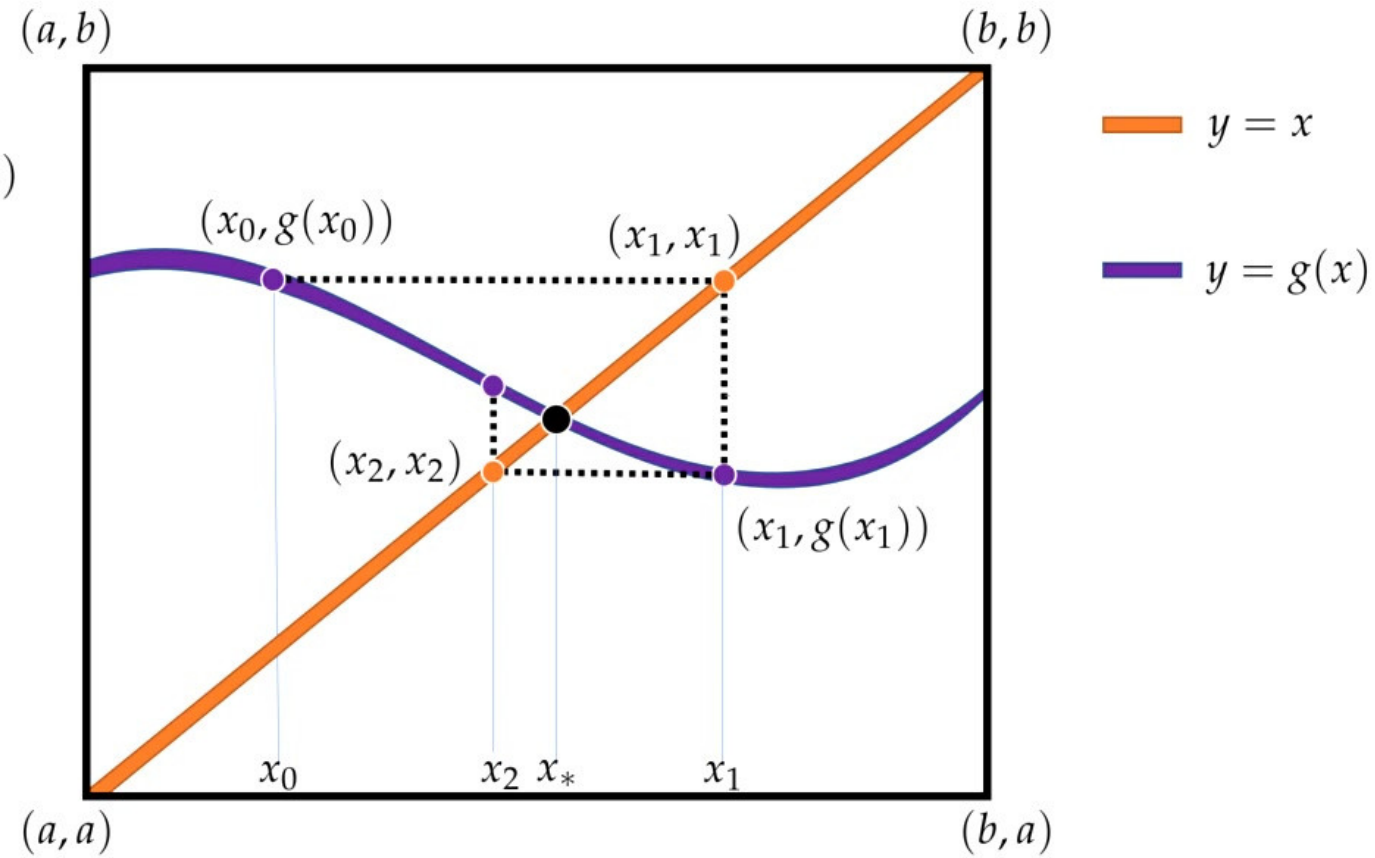

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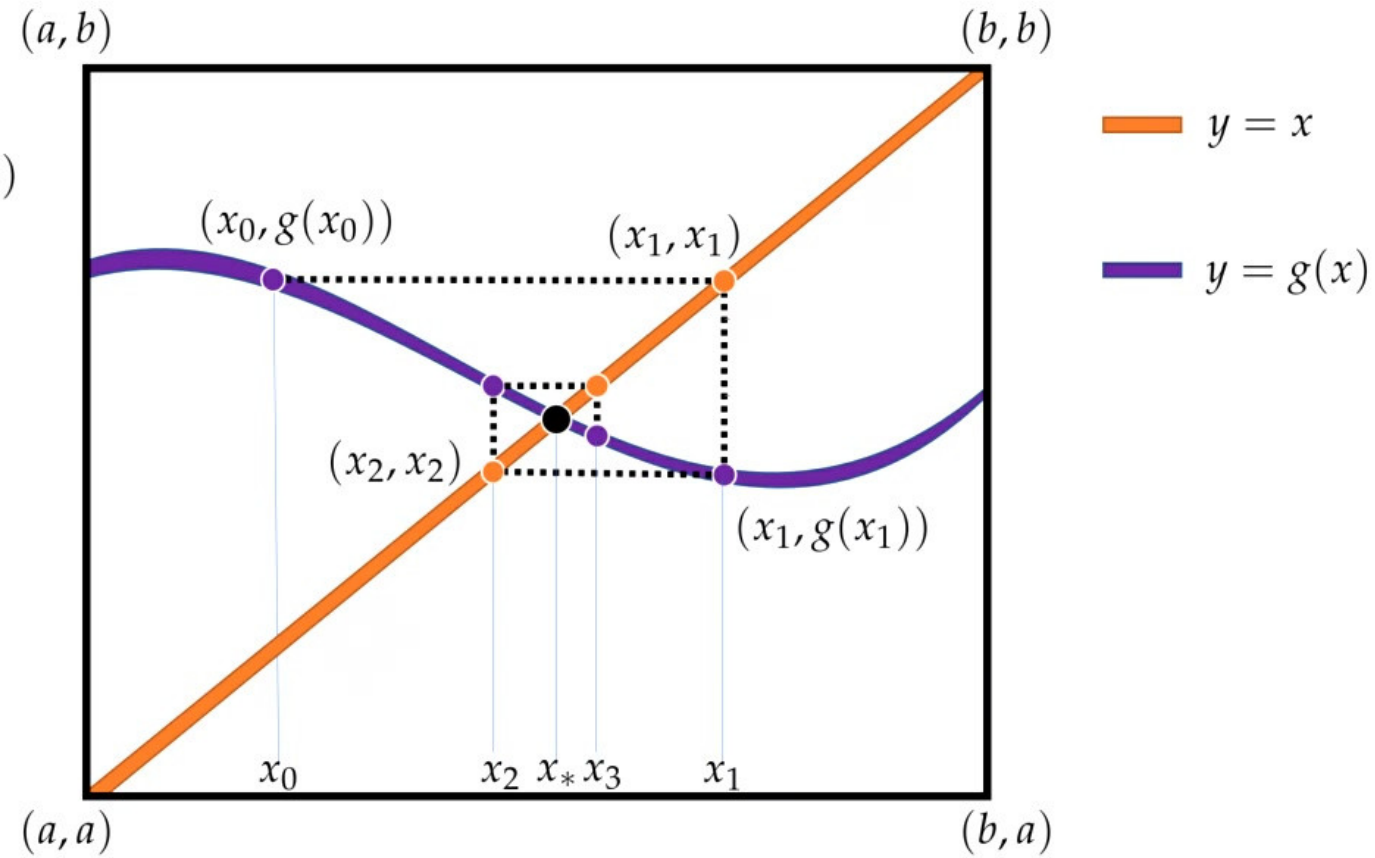

Section 5

The contraction mapping theorem

Let's recall the setting. We have $g : [a, b] \to [a, b]$ with $|g'(x)| < 1$ for $x \in (a, b)$, and we want to find fixed points $x = g(x)$. We know that g has a unique fixed point $x^\star$.

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This algorithm doesn't require derivatives. Can we devise conditions for convergence that don't require derivatives? We'll see this next.

Definition (Contraction)

A function $g : [a, b] \to [a, b]$ is called a *contraction* if there exists a constant $0 \leq \gamma < 1$ such that

$$|g(x) - g(y)| \leq \gamma |x - y|$$

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Any differentiable $g : [a, b] \to [a, b]$ with $|g'(x)| \leq \gamma < 1$ for $x \in (a, b)$ is a contraction! For $x, y \in [a, b]$, by the MVT there exists $c \in (x, y)$ such that

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Not all contractions are differentiable. For example,

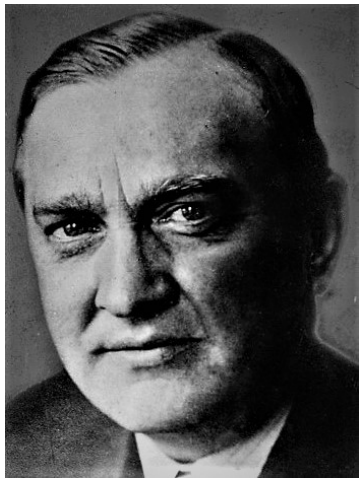
$$g(x) = |x|/2$$

is a contraction with $\gamma = 1/2$, but is not differentiable.

Contraction mapping theorem (1922)

If $g : [a, b] \to [a, b]$ is a contraction, then it has a unique fixed point $x^\star$, and the iteration scheme $x_{i+1} = g(x_i)$ converges at least linearly to $x^\star$ for any $x_0 \in [a, b]$.

Banach proved his theorem on more general *complete metric spaces*.



Stefan Banach, 1892–1945

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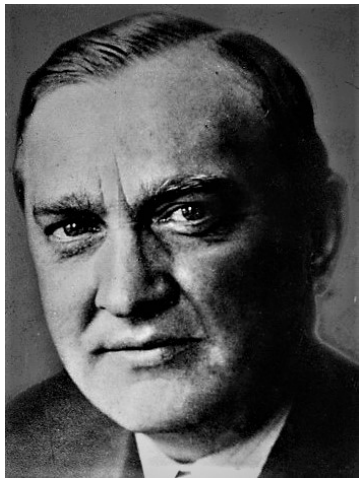
Banach was a Pole who spent his entire academic career in Lwów (now Lviv).



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If $\gamma = 0$ then $g(x) =$ const which is continuous, so assume $\gamma > 0$. Take arbitrary $\varepsilon > 0$ and choose $\delta = \varepsilon/\gamma$. Then if $|x - y| < \delta$, we have

$$|x - y| < \varepsilon/\gamma \implies \gamma|x - y| < \varepsilon,$$

and since $|g(x) - g(y)| \leq \gamma|x - y|$ by assumption, $|g(x) - g(y)| < \varepsilon$.

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We thus know that g must have a fixed point.

Proof.

We now show that the fixed point of g is unique. Suppose p and q are two fixed points of g. Then $g(p) = p$ and $g(q) = q$, so

$$|p - q| = |g(p) - g(q)| \leq \gamma |p - q|$$

and since $\gamma < 1$, this can only be satisfied if $|p - q| = 0$, so $p = q$.

Proof.

We now show convergence for arbitrary $x_0 \in [a, b]$. Recall that $x_i = g(x_{i-1})$ and consider

$$|x_i - x^\star| = |g(x_{i-1}) - g(x^\star)| \leq \gamma |x_{i-1} - x^\star|$$

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Since $\gamma < 1$, $\gamma^i \to 0$, while $|x_0 - x^\star|$ is fixed. Thus

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$$\frac{|x_i - x^\star|}{|x_{i-1} - x^\star|} \leq \gamma,$$

the convergence is at least linear with rate $\gamma < 1$. □

Let's review the conditions for our theorems.

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Let's explore some examples on the edges of these results.

First, let's consider

$$g : [0, 1] \to [0, 1], \quad g(x) = x.$$

This is differentiable but has $|g'(x)| = 1$. Clearly this has an infinite number of fixed points.

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This is differentiable but has $|g'(x)| = 1$. Clearly this has an infinite number of fixed points.

You can have a unique fixed point of a differentiable function without being a contraction. An example is

$$g : [0, \pi] \to [0, 1] \subset [0, \pi], \quad g : x \mapsto \sin x.$$

This has $|g'(x)| < 1$ for $x \in (0, \pi)$, so has a unique fixed point $x^\star = 0$. But it is not a contraction, since $g'(0) = \cos(0) = 1$; there is no $\gamma < 1$ such that $|g'(x)| \leq \gamma$ on $(0, \pi)$. The fixed point iteration converges, but so slowly as to be absolutely useless.

Section 6

Example

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Let's manipulate f to recast the problem as a fixed point problem. There are many ways to do this.

Fixed point iteration A

$$x^2 - x - 1 = 0$$

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$$x^2 - x - 1 = 0 \implies x^2 = x + 1$$

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Fixed point iteration C

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$$x^2 - x - 1 = 0 \implies x = x^2 - 1 =: g_B(x)$$

Fixed point iteration C

$$x^2 - x - 1 = 0 \implies x(x-1) = 1 \implies x = 1/(x-1) =: g_C(x)$$

Comment

This is how the questions for this subject go, but it isn't what rootfinding with fixed point iteration is actually like!

We'll see *generic* ways of transforming a rootfinding problem into a fixed point problem that work for very broad classes of functions.

In other words, the methods actually used don't rely on specific manipulation of the function given.

If we run the fixed point iteration with $x_0 = 1.1$, we get

iteration	$g_A(x) = (x+1)/x$	$g_B(x) = x^2 - 1$	$g_C(x) = 1/(x-1)$
1	1.909091	0.210000	10.00000
2	1.523810	-0.955900	0.111111
3	1.656250	-0.086255	-1.125000
4	1.603774	-0.992560	-0.470588
5	1.623529	-0.014825	-0.680000
6	1.615942	-0.999780	-0.595238
7	1.618834	-0.000439	-0.626866
8	1.617729	-1.000000	-0.614679
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Can we explain this?

Let's check if we can find γ and $[a, b]$ such that $g([a, b]) \subset [a, b]$ and $|g'(x)| \leq \gamma < 1$ on (a, b).

Case A: $g(x) = (x + 1)/x$

Its derivative is $g'(x) = -1/x^2$. On $[a, b] = [1, 2]$ this is increasing, but $g'(1) = -1$. So let's try $[a, b] = [1.1, 2]$. We then have $\gamma = |g'(1.1)| \approx 0.826 < 1$.

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Banach's contraction mapping theorem thus applies.

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Case B: $g(x) = x^2 - 1$

Its derivative is $g'(x) = 2x$. We have $g'(\phi) \approx 3.23 > 1$ and $g'(-\phi^{-1}) \approx -1.23 < -1$. So there can be no interval containing the root that satisfies the criteria.

Let's check if we can find γ and $[a, b]$ such that $g([a, b]) \subset [a, b]$ and $|g'(x)| \leq \gamma < 1$ on (a, b).

Case C: $g(x) = 1/(x - 1)$

Its derivative is $g'(x) = -1/(x - 1)^2$, with $g'(\phi) \approx -2.6 < -1$, and $g'(-\phi^{-1}) \approx -0.38$. Taking $[a, b] = [-0.8, -0.4]$, we have g' is a decreasing function, and $\gamma = |g'(-0.4)| \approx 0.51$.

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On $[-0.8, -0.4]$, g is a decreasing function, so we just need to check the endpoints. We have $g(-0.8) \approx -0.555$ and $g(-0.4) \approx -0.714$, so $g([a, b]) \subset [a, b]$.

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Section 7

Termination criteria

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This reminds us we want a contraction with a small γ: if $\gamma \approx 1$, we will require many iterations to converge.

This is an *a priori* error estimate: we can compute it before ever doing any computations, or choosing x_0. What can we do if we know more?

From the contraction property, we know that

$$|x_i - x_{i-1}| \leq \gamma |x_{i-1} - x_{i-2}|$$

for $i > 2$. Take a fixed $J > i$. We can expand $|x_J - x_i|$ as

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In brackets we have the first few terms of the geometric series, which converges because $\gamma < 1$. Taking the limit $J \to \infty$, so $x_J \to x^\star$, we have

$$|x_i - x^\star| \leq \frac{\gamma^i}{1 - \gamma} |x_1 - x_0|.$$

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In brackets we have the first few terms of the geometric series, which converges because $\gamma < 1$. Taking the limit $J \to \infty$, so $x_J \to x^\star$, we have

$$|x_i - x^\star| \leq \frac{\gamma^i}{1 - \gamma} |x_1 - x_0|.$$

This is an *a posteriori* bound: you have to do some computation to use it.

Section 8

Another example

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Find some $[a, b]$ so that $g(x) = e^{-x}$ has a unique fixed point in $[a, b]$.

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Also note that $g(1) = e^{-1} < 1$, and $g(x)$ is decreasing, so $g : [0, 1] \to [0, 1]$.

We could thus take an interval with $a > 0$ but close and $b = 1$. Choosing $[a, b] = [1/10, 1]$ works fine. (The actual fixed point is $x^\star \approx 0.567143$.)

Continuing with the same example, how many iterations are required to get within 10^{-3} of the fixed point?

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For the *a priori* bound, solving $\gamma^i < 0.001/0.9$ yields $i > 68$. (To achieve a tolerance of 10^{-6}, $i > 137$ is required.)

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Our γ is $e^{-1/10} \approx 0.905$.

For the *a priori* bound, solving $\gamma^i < 0.001/0.9$ yields $i > 68$. (To achieve a tolerance of 10^{-6}, $i > 137$ is required.)

Let's imagine we start with a lucky guess $x_0 = 0.56$. How does the *a posteriori* bound look? In this case $x_1 \approx 0.57120906$, so we have

$$\frac{\gamma^i}{1-\gamma}|0.57120906 - 0.56| < \text{tol},$$

which gives $i > 47$ for $\text{tol} = 10^{-3}$ and $i > 116$ for $\text{tol} = 10^{-6}$.

Section 9

Bonus: accelerating sequence convergence

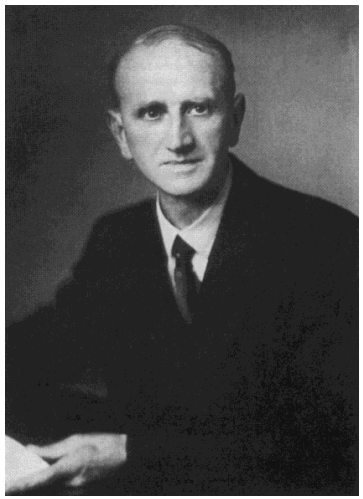
Suppose one has a sequence (x_i) that is linearly converging:

$$\lim_{i \to \infty} \frac{|x_{i+1} - x^\star|}{|x_i - x^\star|} = \mu,$$

with the property that for large enough i,

$$x_i - x^\star, \quad x_{i+1} - x^\star, x_{i+2} - x^\star$$

all have the same sign.



Alexander Aitken FRS FRSL, 1895–1967

Aitken's big idea: use the entries of (x_i) to make a new sequence $(\tilde{x}_i)$ that (hopefully) converges faster!

Assume that the asymptotic limits hold at iterations $i+1$, $i+2$, so that

$$x_{i+1} - x^\star \approx \mu(x_i - x^\star), \quad x_{i+2} - x^\star \approx \mu(x_{i+1} - x^\star).$$

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$$x^\star \approx \frac{(x_i x_{i+2} - x_{i+1}^2)}{x_{i+2} - 2x_{i+1} + x_i}$$

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Aitken thus defines

$$\tilde{x}_i = \frac{\left(x_i x_{i+2} - x_{i+1}^2\right)}{x_{i+2} - 2x_{i+1} + x_i}$$

to yield a new, (hopefully) faster-converging sequence.

Aitken's acceleration is backed up by a theorem.

Aitken's theorem (1926)

Suppose (x_i) is linearly converging with all entries the same sign. Then

$$\lim_{i \to \infty} \frac{\tilde{x}_i - x^\star}{x_i - x^\star} = 0.$$

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Consider Leibniz' formula for π:

$$\pi = 4 \sum_{k=0}^{\infty} \frac{(-1)^k}{2k+1}.$$

Set x_i to be the i^{th} partial sum.

To get π to 10 digits, Leibniz' formula requires about 5 billion terms; Aitken's acceleration $(\tilde{x}_i)$ of it requires about 1400.

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If you apply Aitken acceleration *again*, to yield $(\tilde{\tilde{x}}_i)$, you can get away with only 70 terms!

M4: Constructive Mathematics
Lecture 3: Newton's method

Patrick E. Farrell

University of Oxford

Let's consider rootfinding again:

$$\text{find } x^\star \in \mathbb{R} \text{ such that } f(x^\star) = 0.$$

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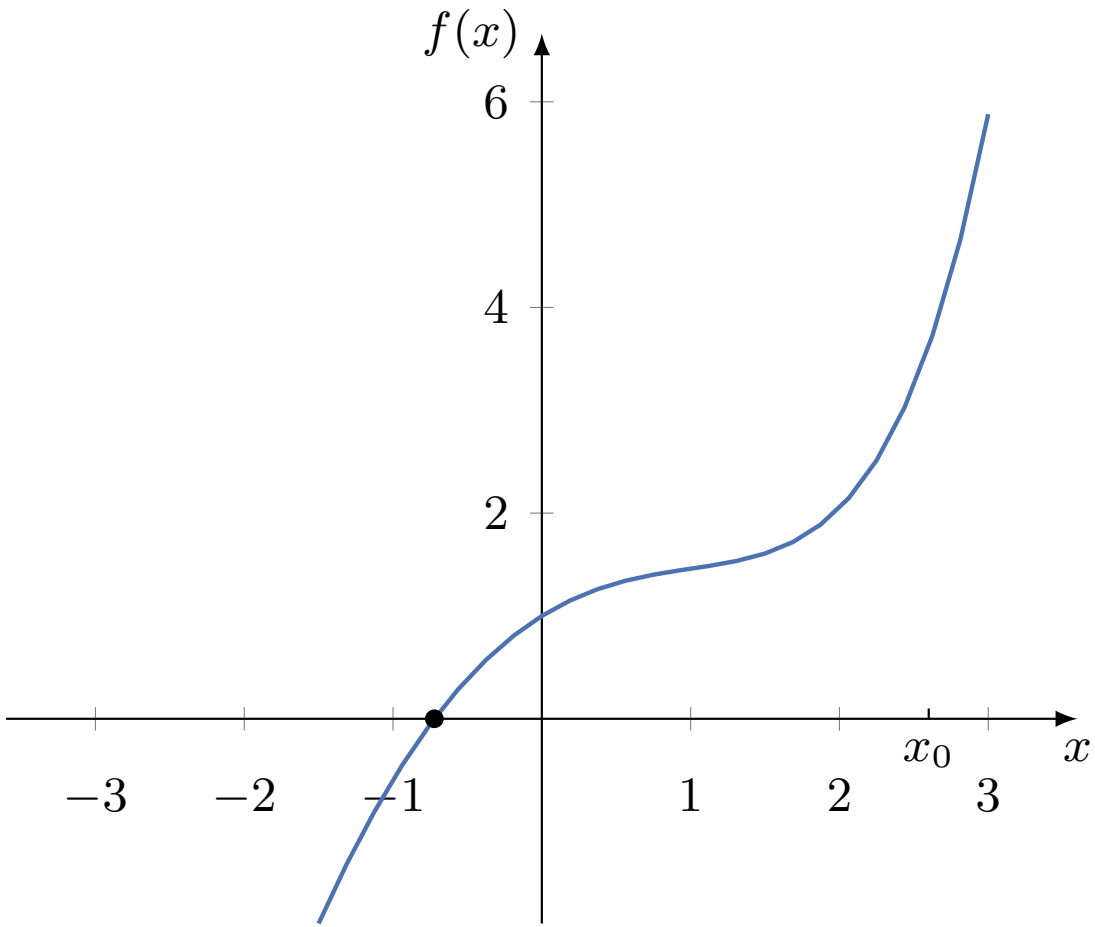
How should we construct $g(x)$ from $f(x)$? One way we've seen is to set

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

but we have no reason to think this is a contraction.

Here is a better way to construct $g(x)$.



Start from an initial x_0.

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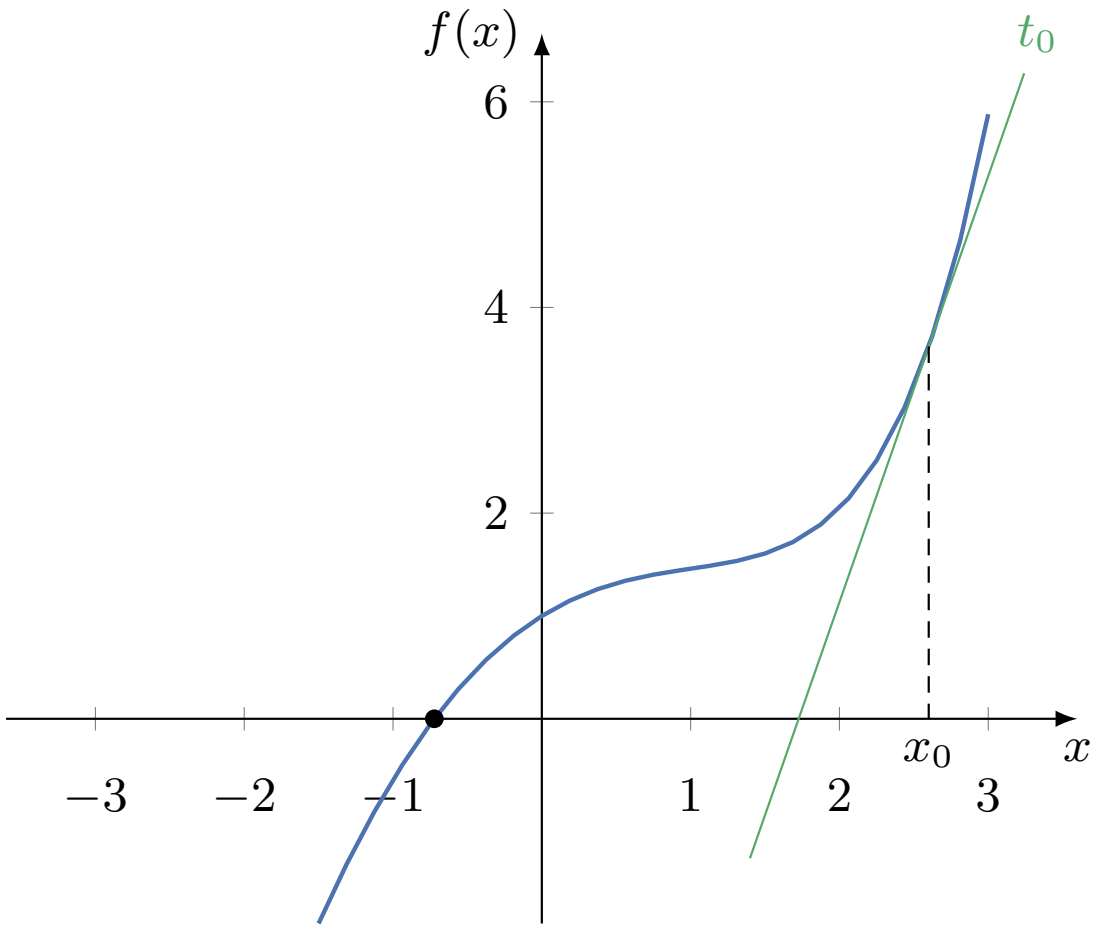

Build a *linear model* of the function.

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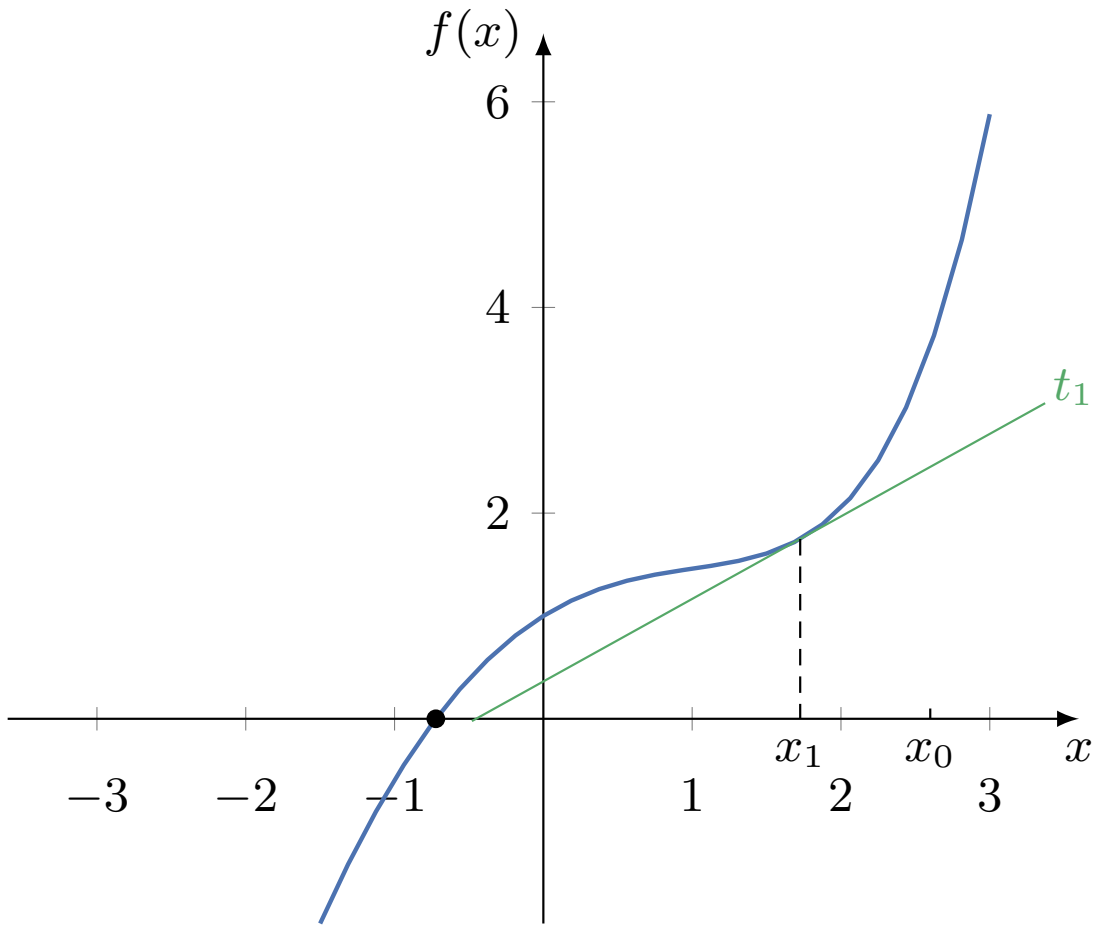

Set x_1 to be the root of the linear model.

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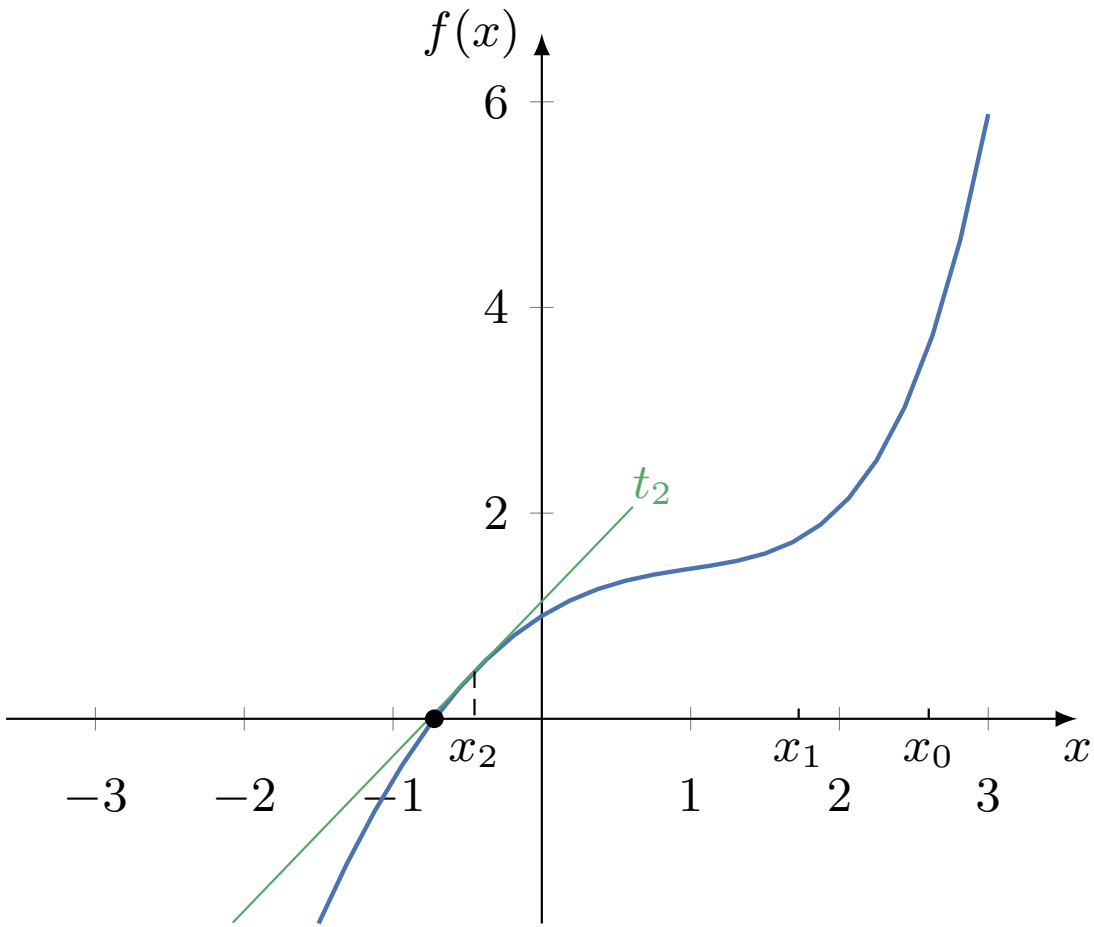

Repeat.

The tangent line joins $(x_i, f(x_i))$ and $(x_{i+1}, 0)$, so we can write its slope as

$$f'(x_i) = \frac{f(x_i) - 0}{x_i - x_{i+1}}$$

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and solving for x_{i+1} yields

$$x_{i+1} = x_i - \left(f'(x_i)\right)^{-1} f(x_i).$$

Newton–Raphson method

$$x_{i+1} = g(x_i) := x_i - \left(f'(x_i)\right)^{-1} f(x_i).$$

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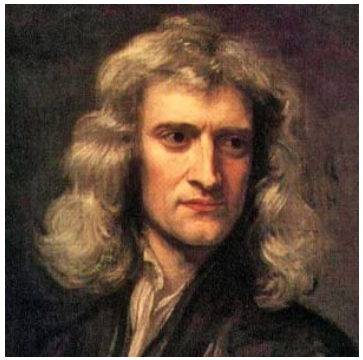

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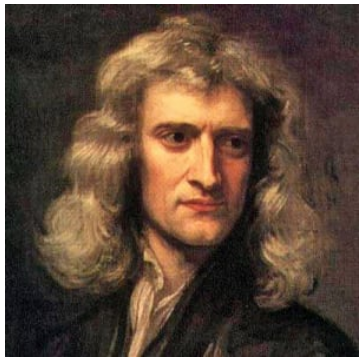

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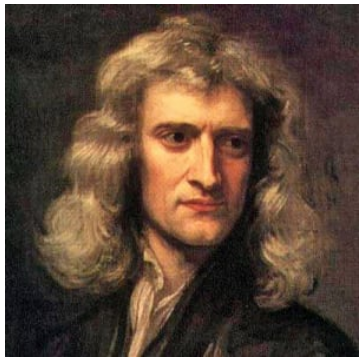

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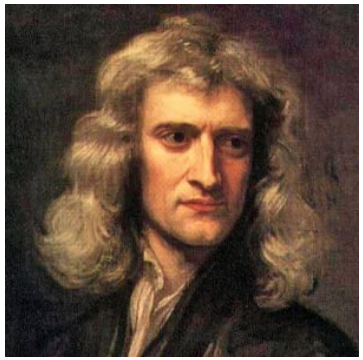

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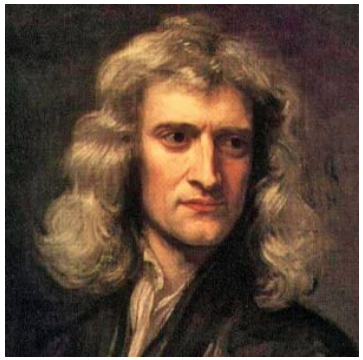

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The extension to computing p-th roots was known to Jamshīd al-Kāshī in Samarkand around 1427.

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Comments:

- ✓ If $f(x_i) = 0$, then $x_{i+1} = x_i$. So roots of f are fixed points of g.

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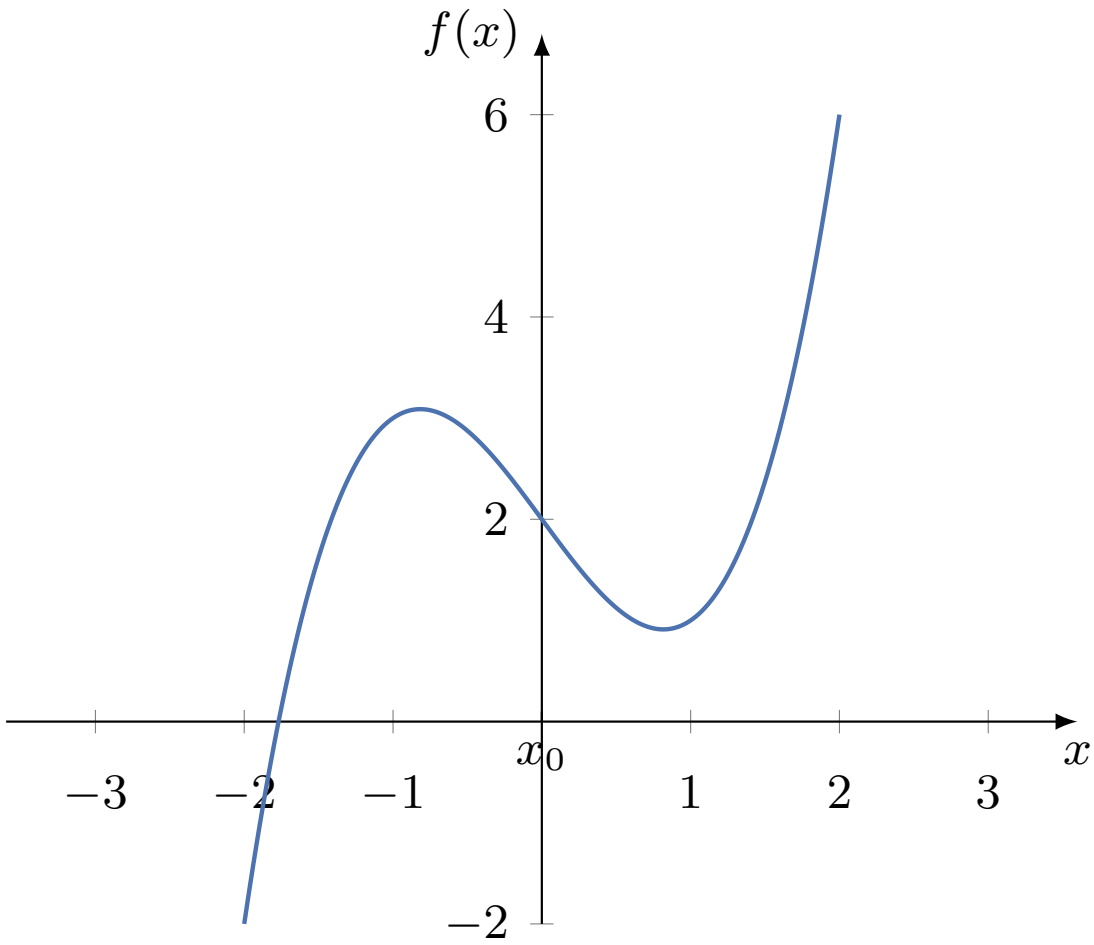
- ✗ If x_0 is far away, the method can diverge or get stuck in a cycle.
- ✓ Newton's method generalises elegantly to higher dimensions.

Consider $f(x) = x^3 - 2x + 2$ with $x_0 = 0$.



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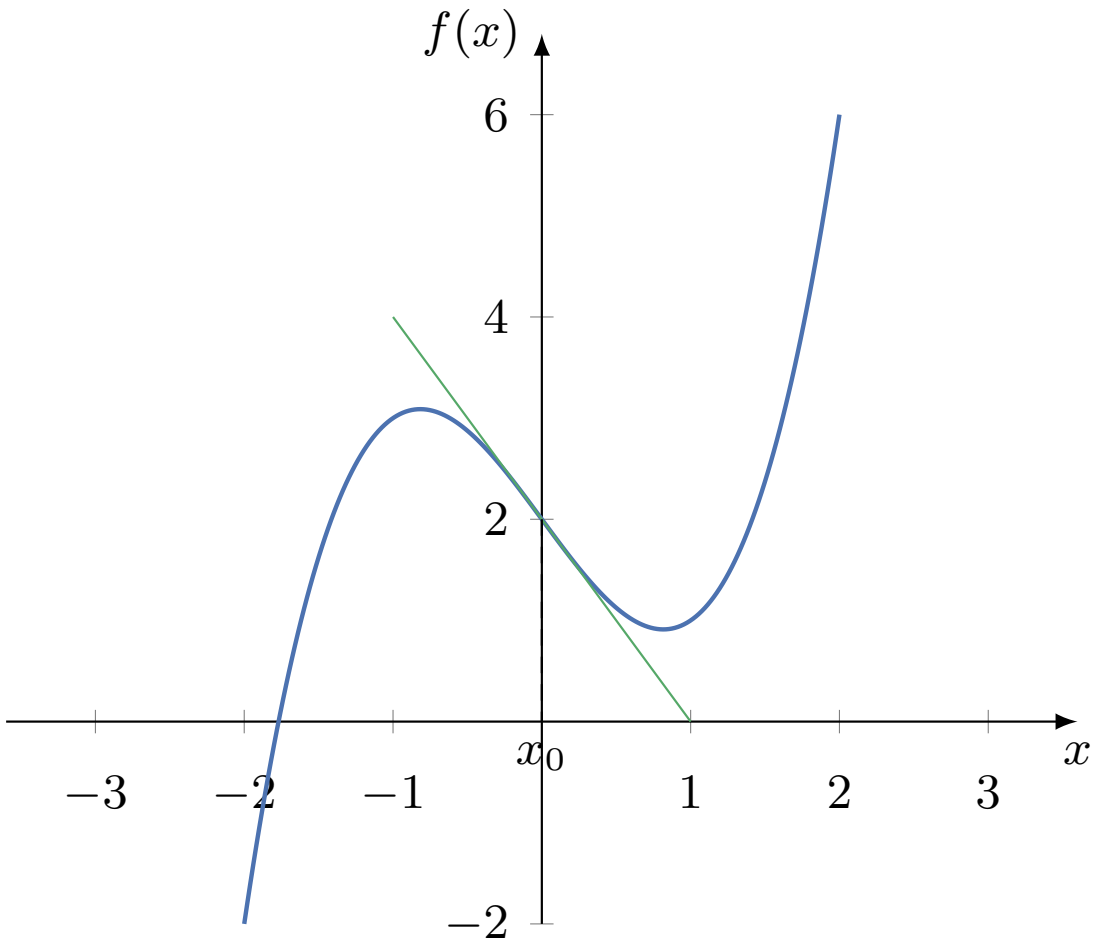

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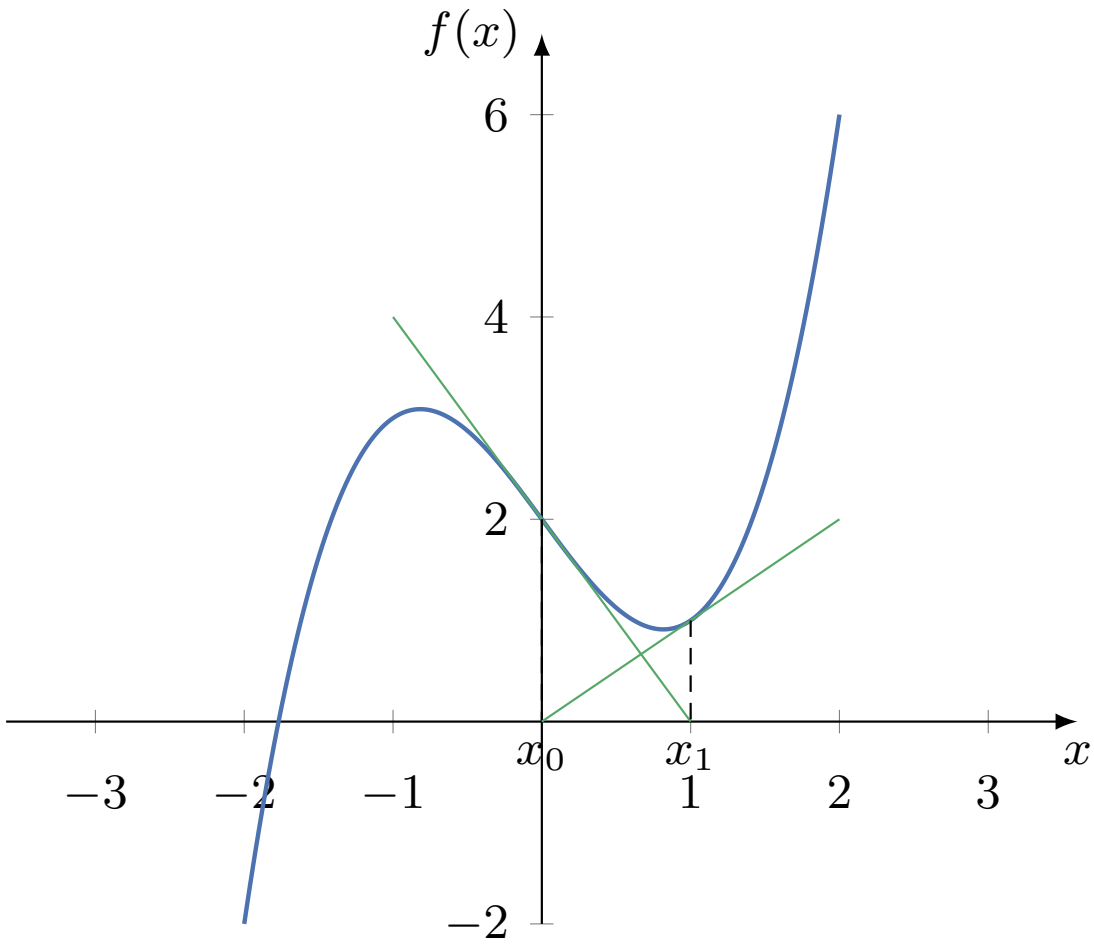

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```
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Iteration  0: x =  2.352836, f(x) = -11.927955
Iteration  1: x = -0.782939, f(x) = 18.906408
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Iteration  3: x = -0.782941, f(x) = 18.906415
...
Iteration  9: x = -0.847671, f(x) = 19.278199
Iteration 10: x =  2.687229, f(x) = -12.596901
Iteration 11: x = -144.9560, f(x) = -3086271.483387
Iteration 12: x = -96.43403, f(x) = -914316.690111
...
Iteration 19: x = -5.622219, f(x) = -167.088934
Iteration 20: x = -4.050607, f(x) = -42.718144
Iteration 21: x = -3.265703, f(x) = -8.235014
Iteration 22: x = -3.023904, f(x) = -0.675602
Iteration 23: x = -3.000221, f(x) = -0.006196
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```

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Iteration  8: x =  2.352281, f(x) = -11.925840
Iteration  9: x = -0.778315, f(x) = 18.878431
Iteration 10: x =  2.332103, f(x) = -11.846922
Iteration 11: x = -0.620547, f(x) = 17.816898
Iteration 12: x =  1.799380, f(x) = -8.442739
Iteration 13: x =  0.804269, f(x) = 2.379590
Iteration 14: x =  0.998101, f(x) = 0.022792
Iteration 15: x =  1.000000, f(x) = 0.000004
Iteration 16: x =  1.000000, f(x) = 0.000000
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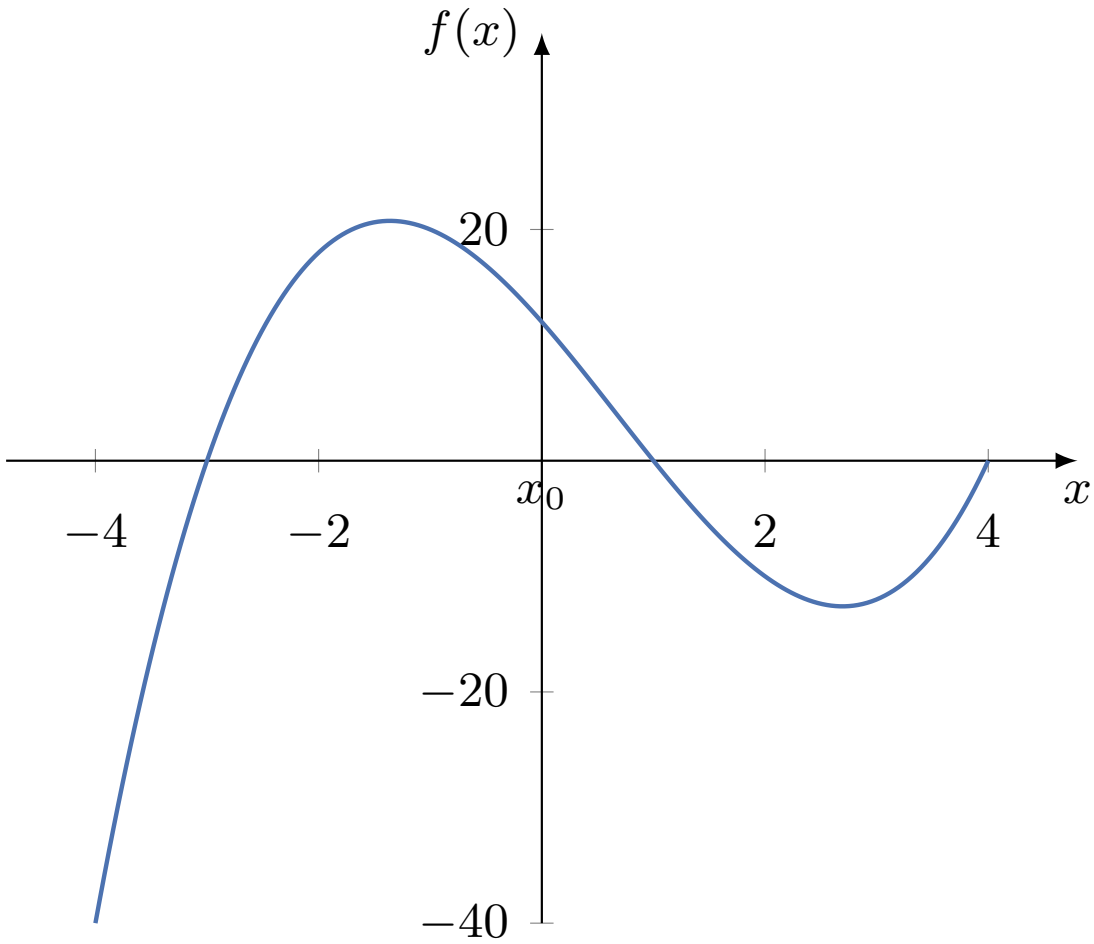
Let's draw the first iterations with $x_0 = 2.352836327$.



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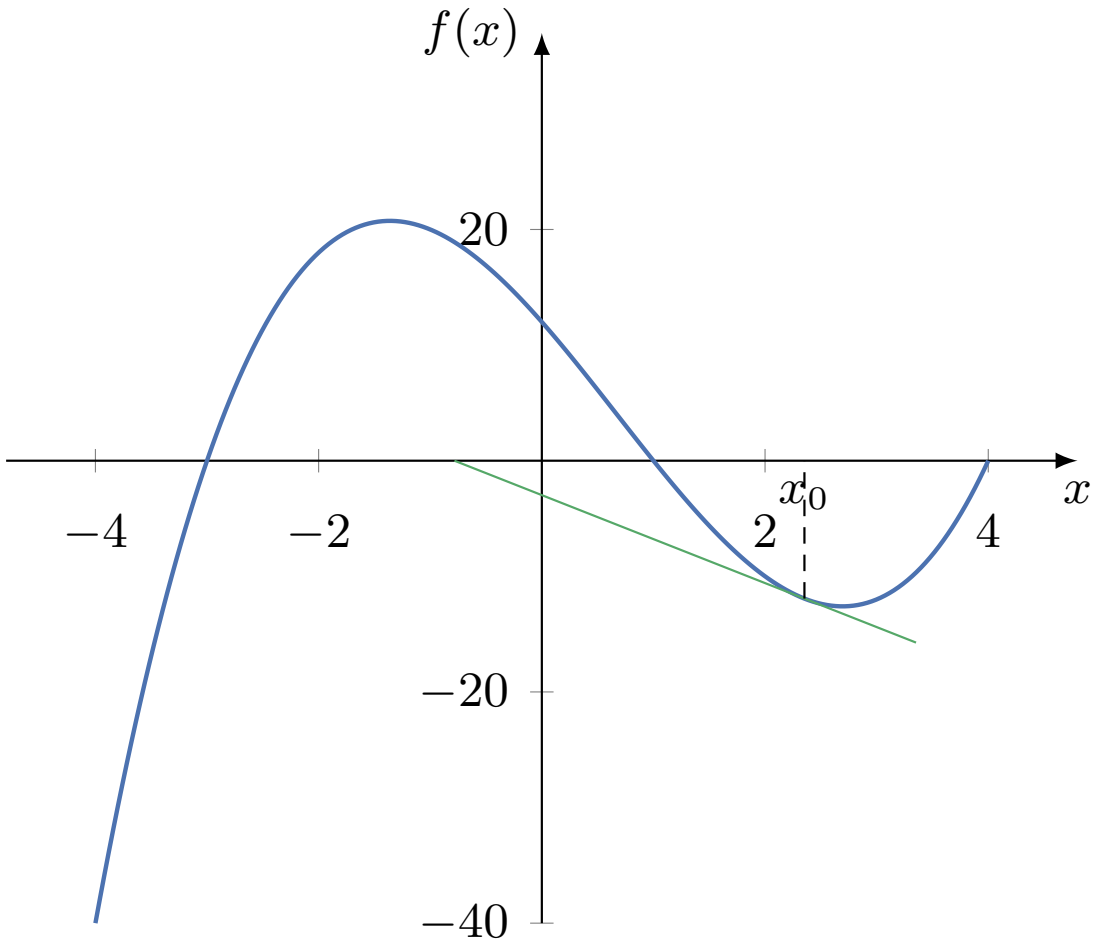

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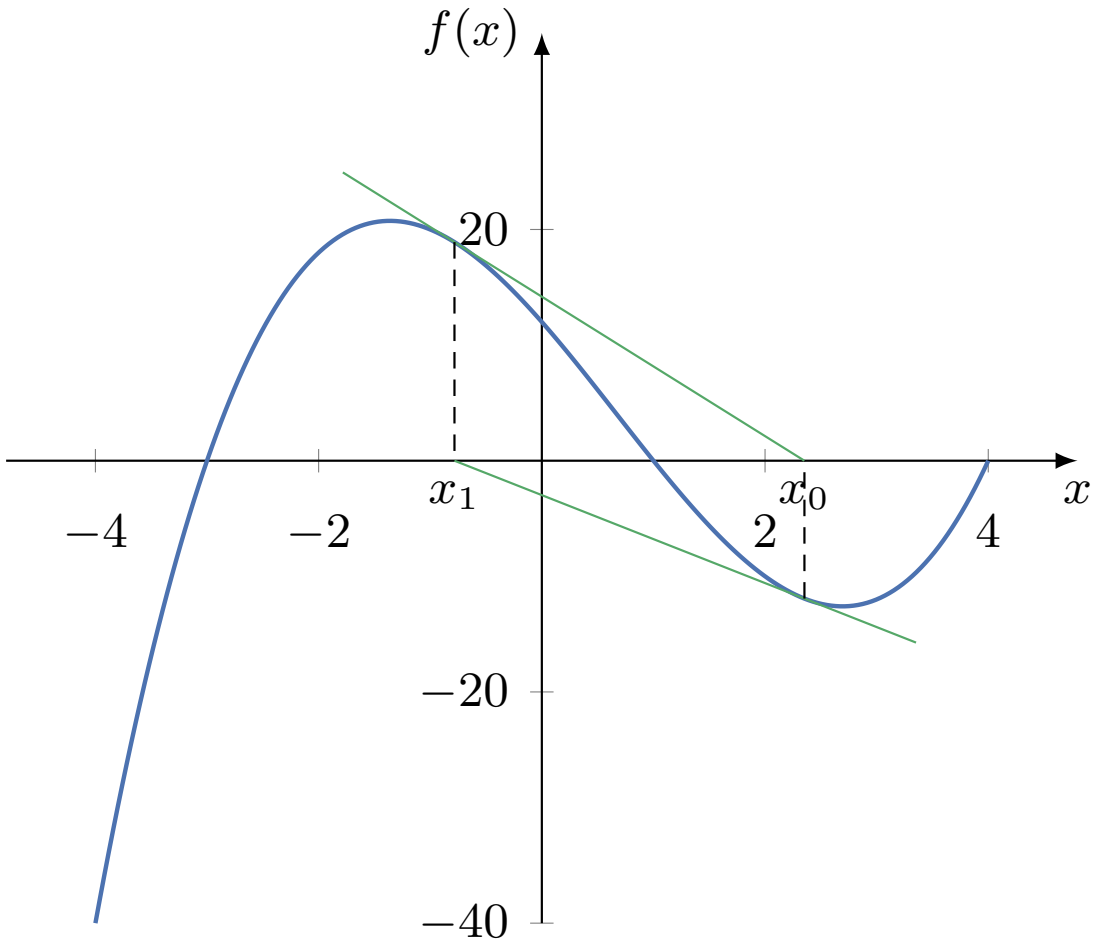

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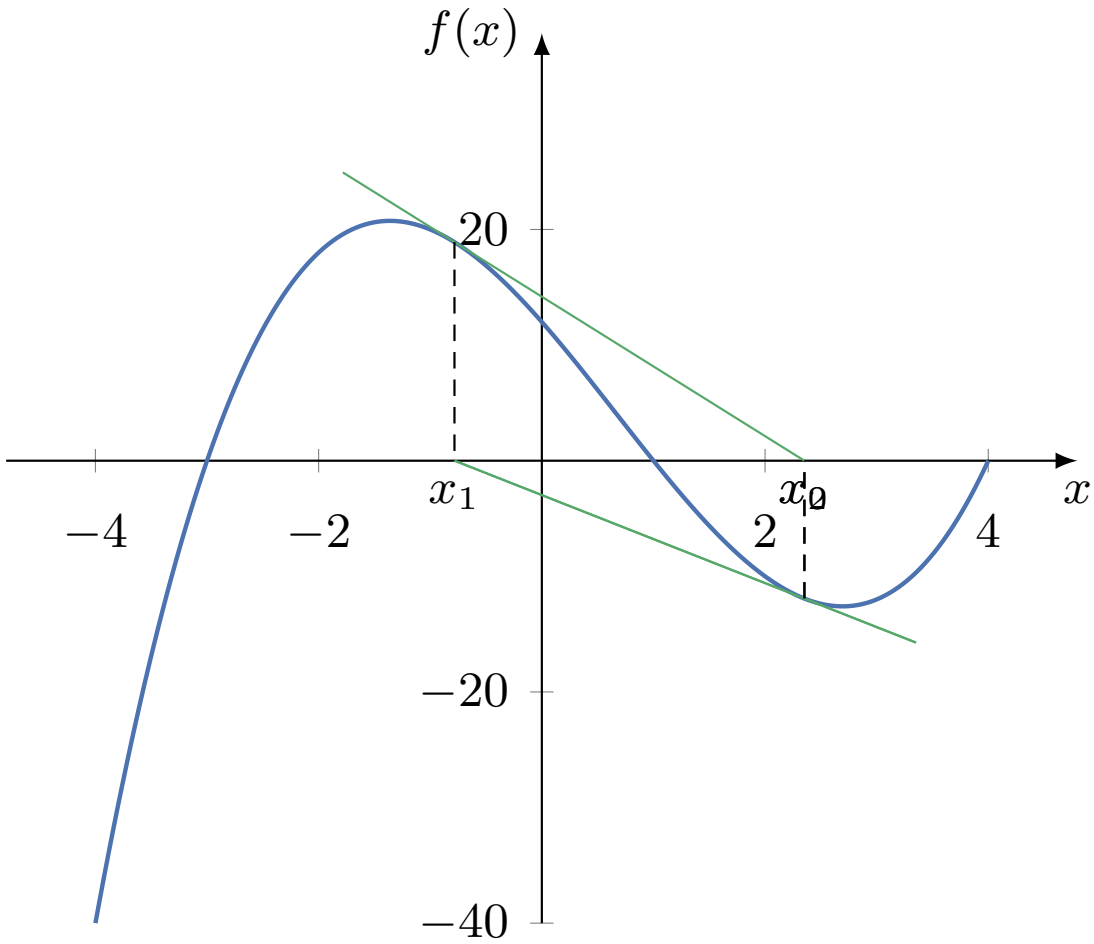

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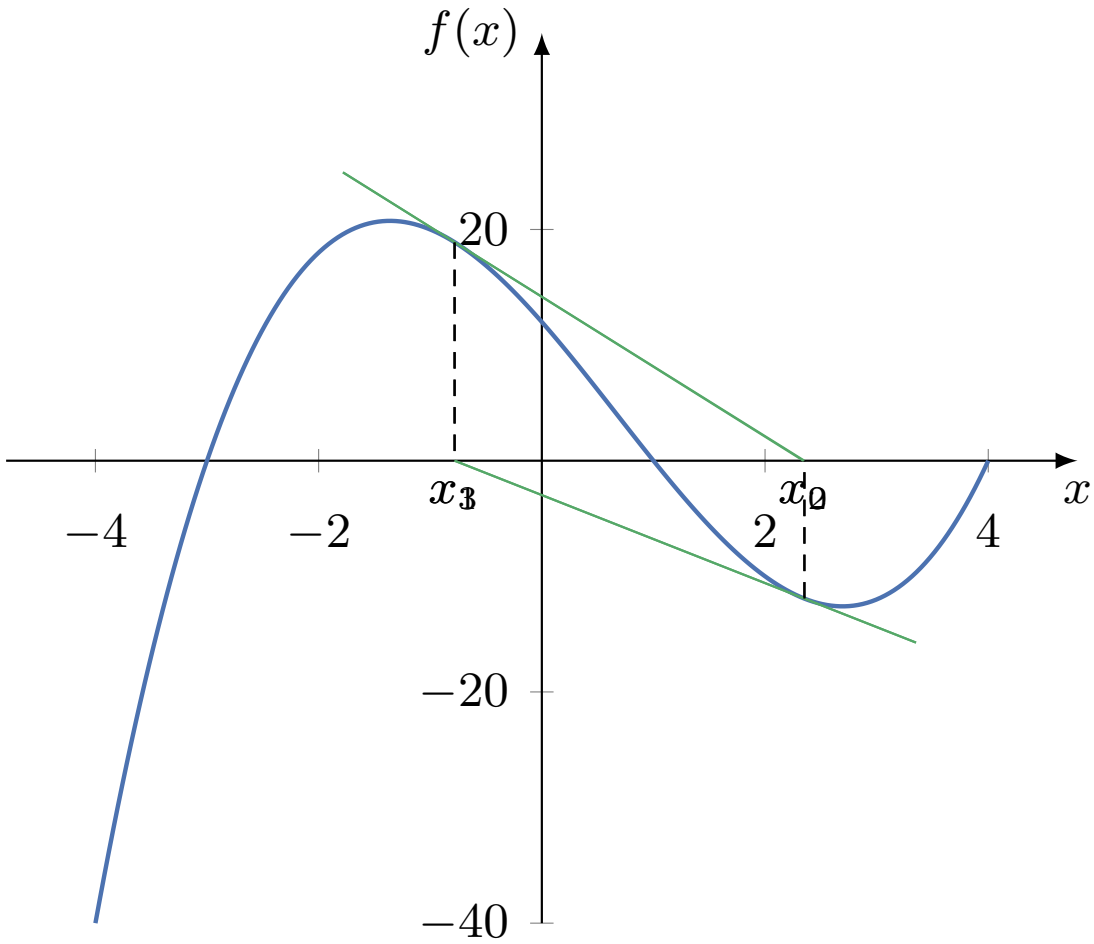

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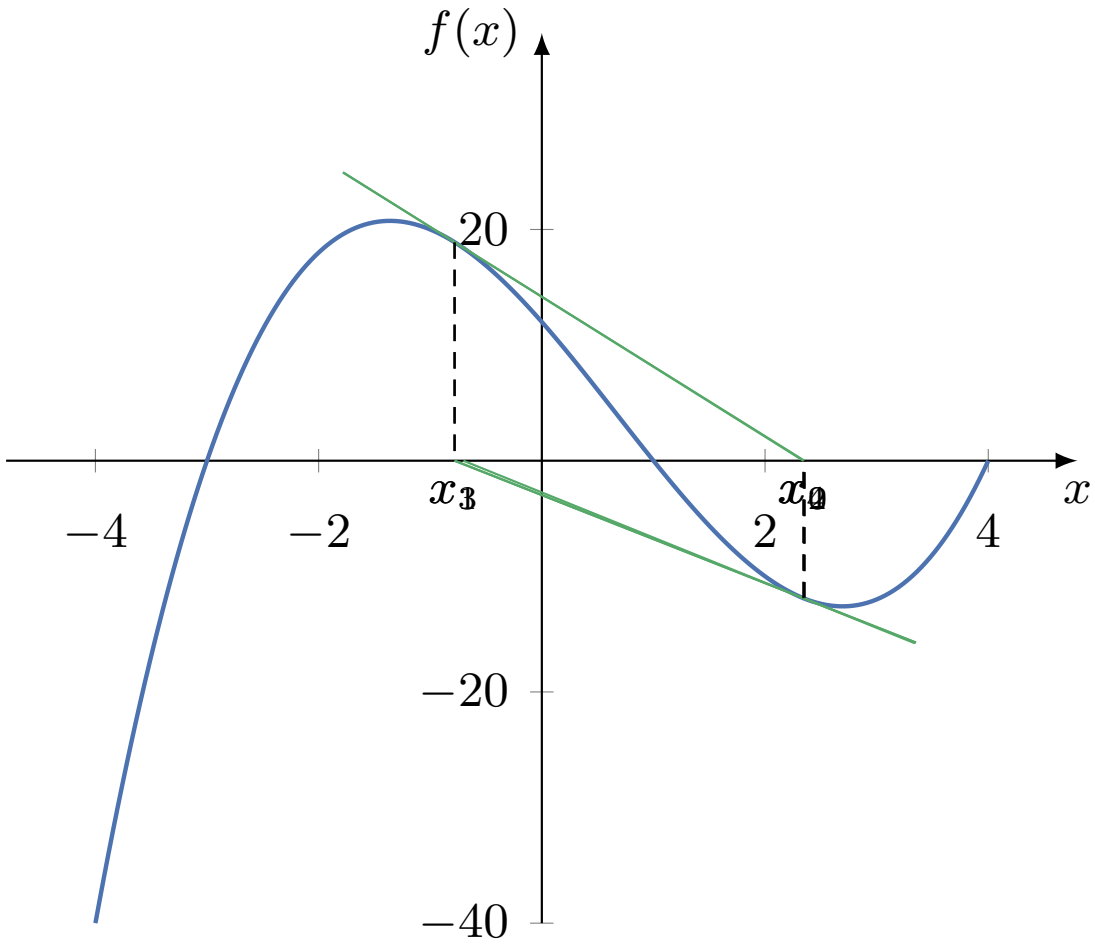

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or in other words that we have linear convergence

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But Newton's method is special: under mild conditions, when x_i is close to $x^\star$ it will satisfy for some $K > 0$

$$\frac{|x_{i+1} - x^\star|}{|x_i - x^\star|^2} \leq K.$$

Recall that we called this *quadratic* convergence.

So why is Newton's method a good idea? Let's talk about general fixed point iterations $x_{i+1} = g(x_i)$ converging to $x^\star$ for a moment.

For a contraction g with contraction factor $\gamma < 1$, we know

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This is much, much faster: roughly speaking, the number of correct digits will *double* at each iteration!

So how can Newton achieve this quadratic convergence?

Recall the Taylor expansion of g around some point a:

$$g(x_i) = g(a) + (x_i - a)g'(a) + \frac{1}{2}(x_i - a)^2 g''(\zeta_i), \quad \text{some } \zeta_i \in (x_i, a).$$

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If $f'(x^\star) = 0$, we have a multiple root, and we have to take the limit $x \to x^\star$ and use L'Hôpital's rule to evaluate the fraction.

Take-home message

Newton's method converges quadratically to isolated roots.

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If the root is not isolated, then one generally expects linear convergence, with the exact rate depending on details. For example, on the problem sheets you will prove that if

$$f'(x^\star) = 0, f''(x^\star) \neq 0$$

then one expects linear convergence with rate $1/2$.

Let's take an example. Let's look for the fixed point of $x = \cos x$. We tried this with bisection and it was slow.

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Let's do an exam question. Consider the question from 2017, Paper IV, Q7 (b):

> *The function*
>
> $$p(x) = 27x^3 - 27x^2 + 4$$
>
> *has a root* $\alpha = 2/3$.
>
> *Show that Newton's method to compute approximations to this root, with starting guess* x_0, *can be written as the iteration*
>
> $$x_{k+1} = g(x_k),$$
>
> *where you should find* g *explicitly. Prove or disprove that the sequence generated will converge to* α *for any* $x_0 \in [1/3, 1]$.

We write

$$\begin{aligned} g(x) &= x - \frac{p(x)}{p'(x)} \\ &= x - \frac{27x^3 - 27x^2 + 4}{81x^2 - 54x} \end{aligned}$$

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To check whether the Newton sequence will converge, we investigate the conditions of Banach's contraction mapping theorem.

Let's check the conditions. We compute

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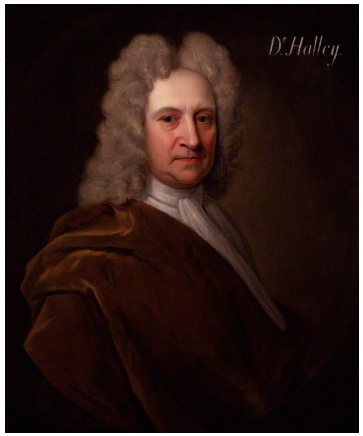
So the conditions of Banach's contraction mapping theorem are satisfied.

There are other fixed-point iterations for rootfinding.

Halley's method (1694)

$$x_{i+1} = g(x_i) := x_i - \frac{2f(x_i)f'(x_i)}{2[f'(x_i)]^2 - f(x_i)f''(x_i)}.$$



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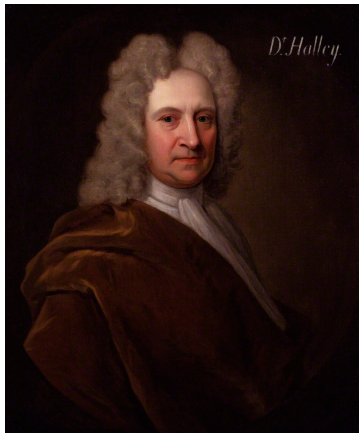

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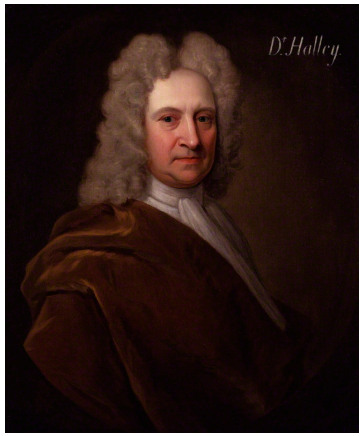

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In a letter in 1712, Taylor wrote

> *While I was thinking of these things, I fell into a general method of applying Dr. Halley's Extraction of roots to all Problems ...And it is comprehended in this Theorem ....*

The theorem he proved was Taylor's theorem!

Brook Taylor FRS, 1685–1731

Section 2

Bonus: the secant iteration

Halley's method uses more derivatives to get faster convergence.

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The *secant* iteration makes the converse trade: no derivative evaluations, for (slightly) slower convergence.

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$$f'(x_i) \approx \frac{f(x_i) - f(x_{i-1})}{x_i - x_{i-1}}$$

with some previous data x_{i-1}.

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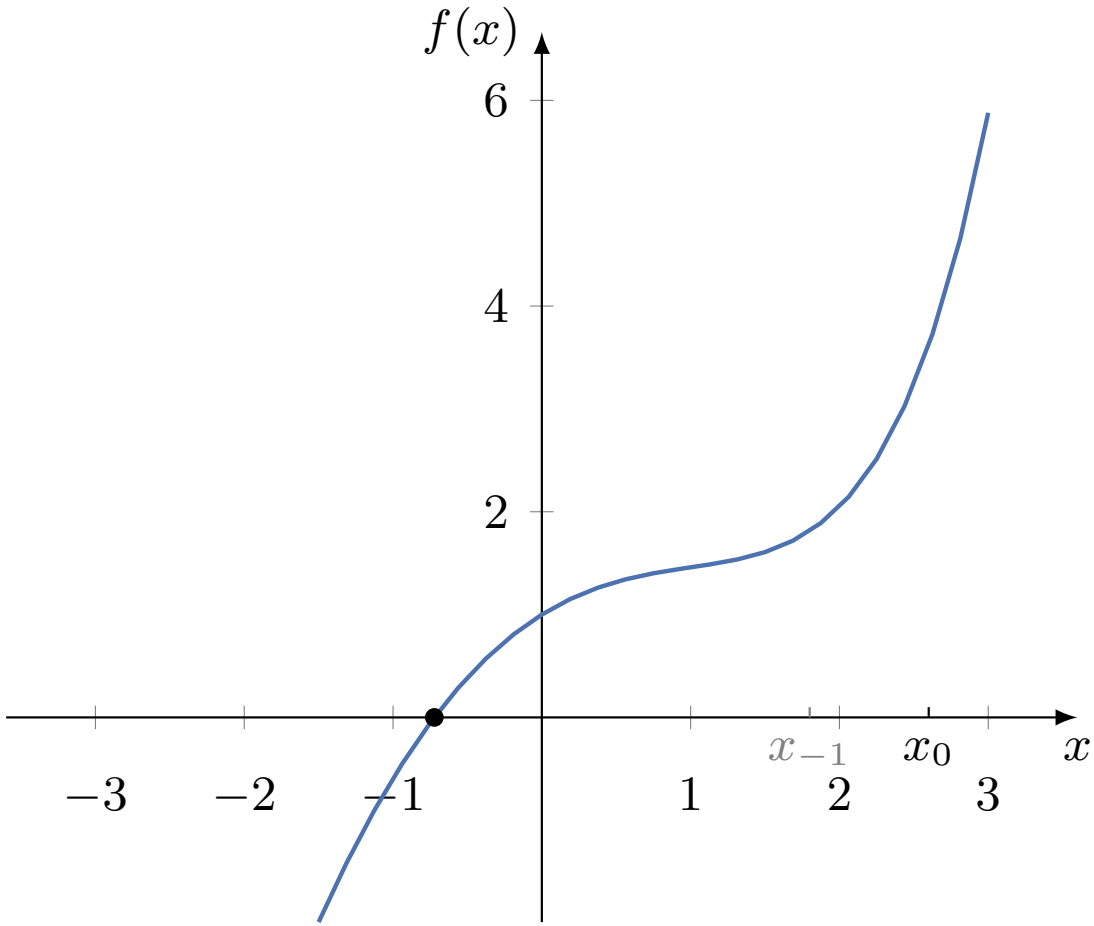
Newton invented the secant method around the same time, but never published it.

Both the ancient Egyptians and Babylonians used the secant method around 1800 BCE to solve equations like

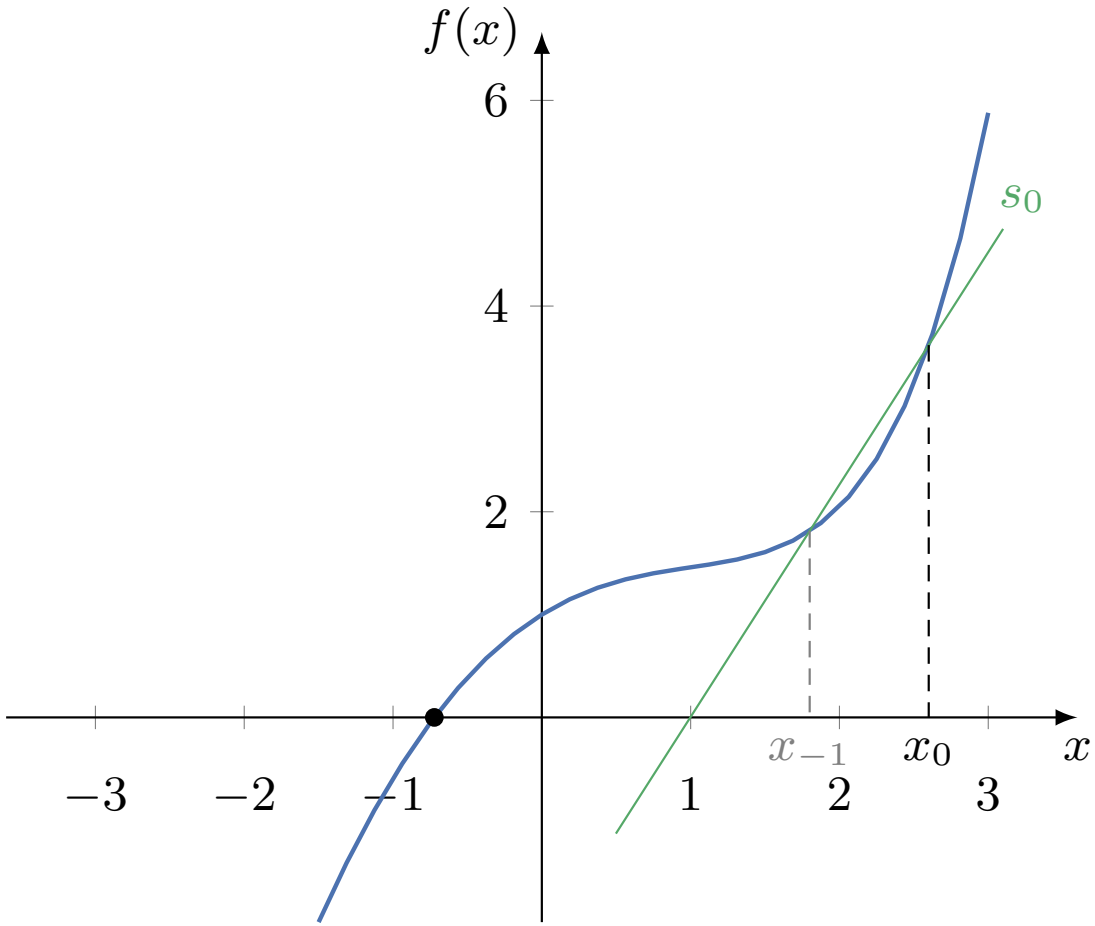
$$ax + b = c$$

since they didn't know how to move terms from one side to another!

The secant iteration.



The secant iteration.



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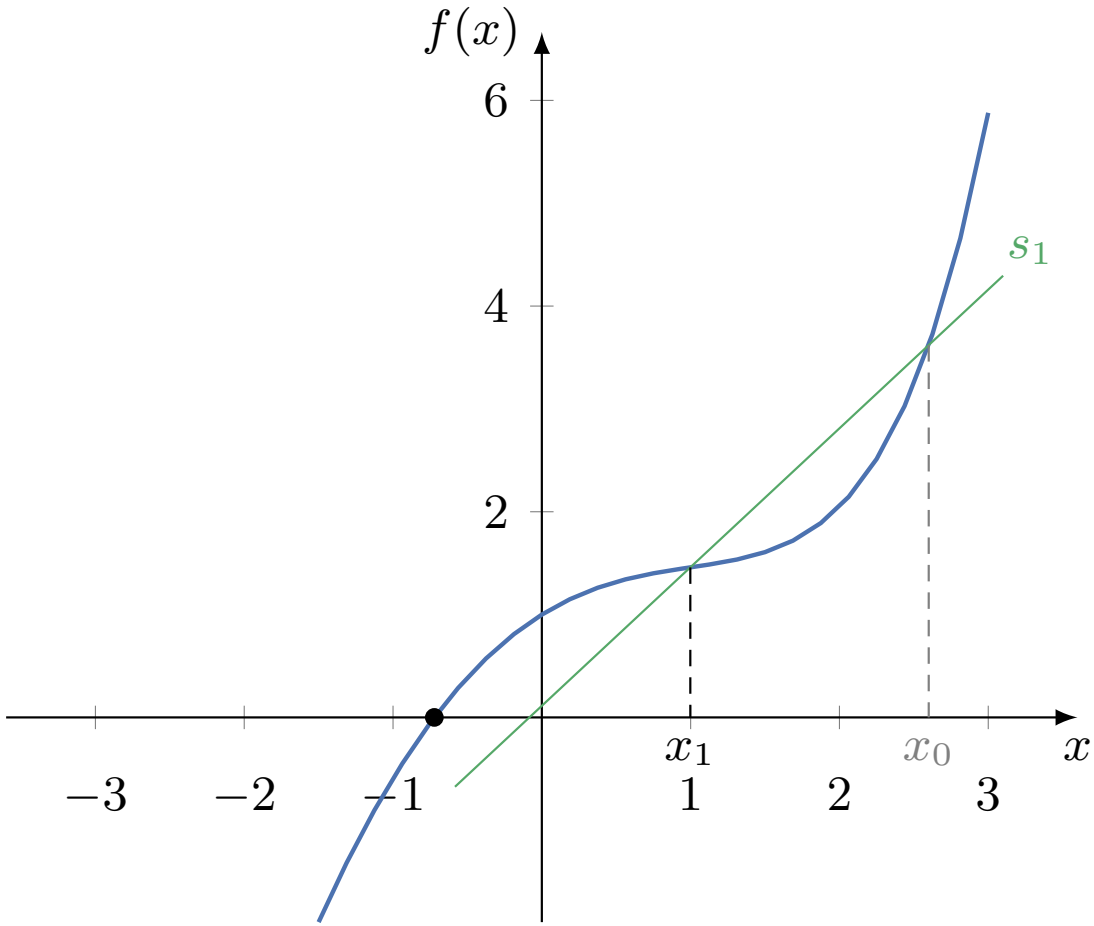

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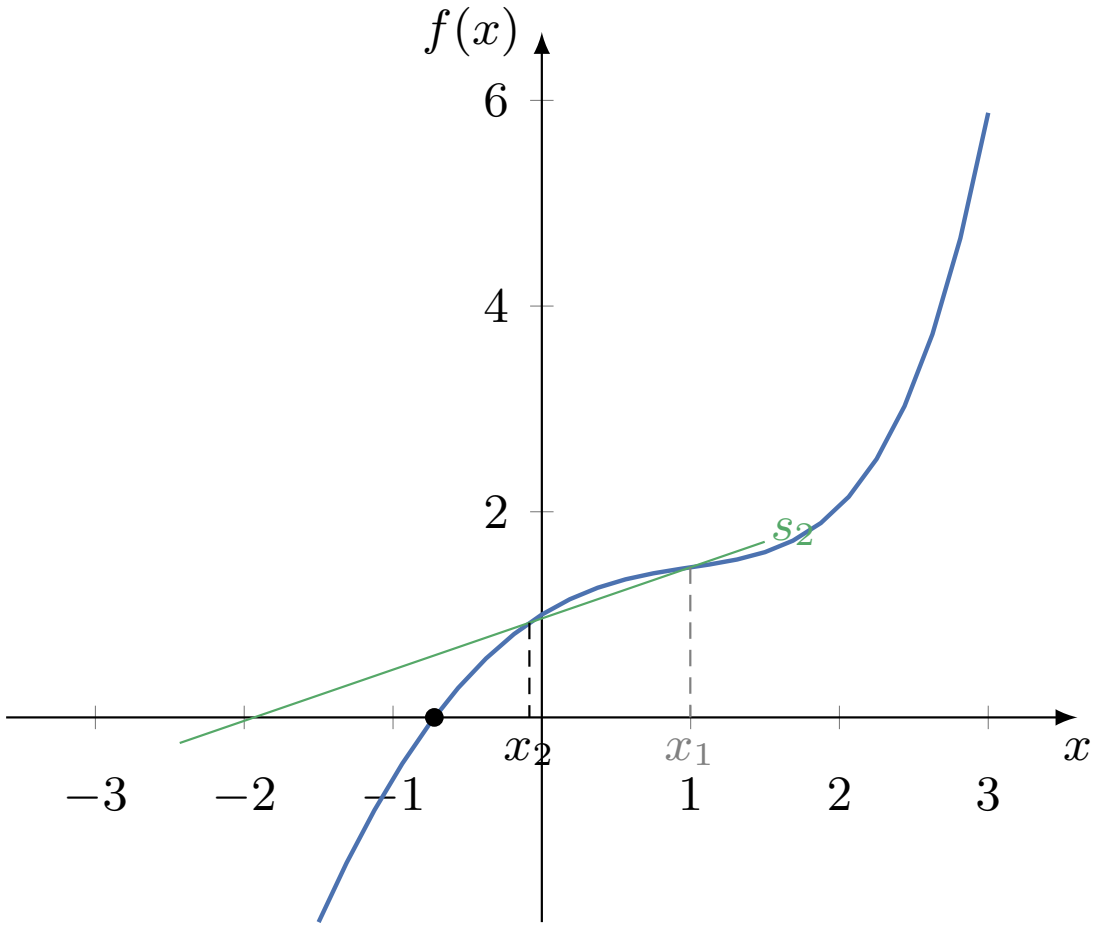

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Comments on the secant method:

- ✗ The method requires more information to start, and depends sensitively on it.
- ✓ In principle the method can be applied to nondifferentiable functions.
- ▶ The generalisation to higher dimensions is different—leading to the quasi-Newton family of methods.

Section 3

Bonus: Aitken acceleration of fixed-point iterations

Suppose our fixed-point iteration

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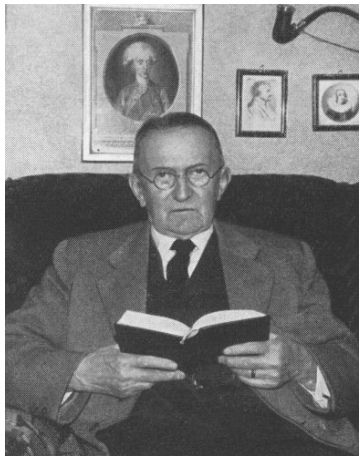
The acceleration only goes one way: we don't re-use the accelerated values in the fixed-point iteration itself.

Steffensen's idea

Do two steps of fixed-point iteration, apply Aitken acceleration, then re-start the fixed-point iteration from there.

This *interleaves* the fixed-point iteration and acceleration.



Johan Frederik Steffensen, 1873–1961

Assume $g : [a, b] \to [a, b]$, and $x_0 \in [a, b]$.

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---

function steffensen(g, x_0, tol)

  $x \leftarrow x_0$

  while $|g(x) - x| > \text{tol}$ **do**

    $x_0 \leftarrow x$

    $x_1 \leftarrow g(x_0)$

    $x_2 \leftarrow g(x_1)$

    $x \leftarrow \left(x_0 x_2 - x_1^2\right) / (x_2 - 2x_1 + x_0)$

  end while

  return $g(x)$

end function

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If you organise the code properly, this requires two evaluations of g per iteration.

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  return $g(x)$
end function

If you organise the code properly, this requires two evaluations of g per iteration.

Does this really help?

Yes, it does, under certain conditions:

Steffensen's theorem (1933)

Suppose that $g(x)$ has a fixed point $x^\star$ with $g'(x^\star) \neq 1$. If there exists $\delta > 0$ such that $g \in C^3([x^\star - \delta, x^\star + \delta], \mathbb{R})$, then Steffensen's method gives quadratic convergence for any $x_0 \in [x^\star - \delta, x^\star + \delta]$.

This can achieve quadratic convergence, without derivatives!

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Fixed-point iteration requires 37 evaluations of g to get ϕ to 16 digits. Steffensen's method requires only 8!

Let's apply Newton's method to

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```
>> newton(@(x) (x-1)^2, @(x) 2*x - 2, 2, 1e-4)
Iteration  0: x = 2.000000, f(x) = 1.000000
Iteration  1: x = 1.500000, f(x) = 0.250000
Iteration  2: x = 1.250000, f(x) = 0.062500
Iteration  3: x = 1.125000, f(x) = 0.015625
Iteration  4: x = 1.062500, f(x) = 0.003906
Iteration  5: x = 1.031250, f(x) = 0.000977
Iteration  6: x = 1.015625, f(x) = 0.000244
Iteration  7: x = 1.007812, f(x) = 0.000061
```

Converging linearly, you say?

```
>> steffensen(@(x) x - (x-1)^2/(2*x - 2), 2, 1e-4)
Iteration  0: fixed point = 2.000000
Iteration  2: fixed point = 1.000000
```

Steffensen's method gets the answer exact to 16 digits in 2 iterations.

Section 4

Rootfinding for polynomials

We have seen general rootfinding methods that apply to many different kinds of functions.

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Philosophical remark

When designing algorithms, we should always ask: have we used every piece of knowledge we have about the problem?

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When designing algorithms, we should always ask: have we used every piece of knowledge we have about the problem?

For example, if we restrict ourselves to rootfinding for *polynomials*, can we make our algorithms better? The answer is yes.

Section 5

Horner's method

In the literature, Horner's method refers to two different things:

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The evaluation scheme was known in medieval times to Qín Jiǔsháo (c. 1202–1261) and Sharaf al-Dīn al-Ṭūsī (c. 1135-1213), and later to Newton and Lagrange.

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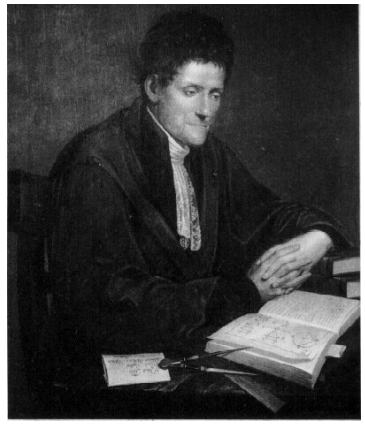
It's not clear that Horner, a schoolmaster in Bath, even invented the latter method that now bears his name. He was beaten to it by Paolo Ruffini in 1804 and Theophilus Holdred, a London watchmaker, in 1820. The method was published again by Horner in 1830.



Paolo Ruffini, 1765–1822

Let's consider Horner's two methods in order. Suppose we have a polynomial

$$p(x) = a_0 + a_1 x + a_2 x^2 + \cdots + a_n x^n$$

with n large, e.g. $n = 10,000$. How should we evaluate $p(r)$ for $r \in \mathbb{R}$?

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with n large, e.g. $n = 10,000$. How should we evaluate $p(r)$ for $r \in \mathbb{R}$?

One way would be to evaluate all the terms in the sum separately, and add them up. This would require n additions and

$$0 + 1 + 2 + \cdots + n = \frac{n^2 + n}{2}$$

multiplications. Scaling like n^2 is bad!

Instead, a better way is to write

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This shares the evaluations of powers of x. It only requires n multiplications and n additions. Much faster!

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We then have $b_0 = p(r)$.

There's more to it than this, however.

Theorem

Define the polynomial

$$Q(x) := b_n x^{n-1} + b_{n-1} x^{n-2} + \cdots + b_2 x + b_1.$$

Then

$$p(x) = (x - r)Q(x) + b_0.$$

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so in particular

$$p'(r) = Q(r).$$

Proof.

Recall that $p(x) = a_0 + \cdots + a_n x^n$, $b_n = a_n$, and $b_i = a_i + b_{i+1} r$.

Expand

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In the context of Newton's method applied to p, we have

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$$x_{i+1} = x_i - \frac{p(x_i)}{p'(x_i)} = x_i - \frac{p(x_i)}{Q(x_i)}.$$

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$b \leftarrow a_n x + a_{n-1}$ # Horner eval for p

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if $|b| <$ tol **then** # success

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        if $|b| <$ tol **then** # success
            return x
        end if
        $x \leftarrow x - b/c$ # Newton update
    end for
end function

We can summarise with the following useful notation:

Definition (Big $\mathcal{O}$ notation)

For $g(n) > 0$, we say

$$f(n) = \mathcal{O}(g(n)) \text{ as } n \to \infty$$

if there exists $M > 0$ and $n_0 \in \mathbb{N}$ such that

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The number of operations to evaluate a degree-n polynomial is:

- $\mathcal{O}(n^2)$ for the naïve way, but
- $\mathcal{O}(n)$ for Horner's evaluation scheme.

This is much, much better at high n!

In fact, Horner's scheme for evaluation has a nice optimality property:

Theorem

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If you know you'll evaluate a polynomial many times on different inputs, it is possible to preprocess the polynomial into a representation that requires fewer operations (trading offline work for online work).

Section 6

More philosophical remarks

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Philosophical remark

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$$a_0 + a_1 x + \cdots + a_n x^n = a_0 + x\left(a_1 + x\left(a_2 + \cdots + x\left(a_{n-1} + x a_n\right)\cdots\right)\right).$$

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Algorithmic advances sometimes come by deriving an equivalent expression with better properties.

Think back to our list of questions we ask about algorithms:

- Does the algorithm terminate?
- Does the algorithm give the correct answer?
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Every computer nowadays has multiple processing units. (My phone has 8.) Can we use them?

Here's another equivalent expression with different properties:

$$\begin{aligned} &a_0 + a_1 x + \cdots + a_n x^n \\ &\quad = \left(a_0 + a_2 x^2 + a_4 x^4 + \cdots\right) + \left(a_1 x + a_3 x^3 + a_5 x^5 + \cdots\right) \end{aligned}$$

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which we can evaluate in parallel with two independent runs of Horner's method.

More generally, if you have enough terms, you can break p up into $k+1$ polynomials $\{p_j\}_{j=0}^{k}$, each taking the monomial term x^i if

$$i \bmod (k+1) = j.$$

Section 7

Bonus: finding all roots of a polynomial

Horner's scheme is just a specialised variant of Newton's method. It finds roots one at a time.

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Once you have found a root $x^\star$ of $p_0(x)$, you can construct

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and apply the scheme again to p_1. Iterating in this way one can find all real roots, if you can construct good initial guesses.

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Can we find them all at once, without fussing over guesses?

It turns out that we have *very* fast and powerful algorithms for computing the eigenvalues of matrices:

$$\text{for } A \in \mathbb{R}^{n \times n}, \text{ find all } \lambda_i, v_i \text{ s.t. } Av_i = \lambda v_i, \|v_i\|^2 = 1.$$

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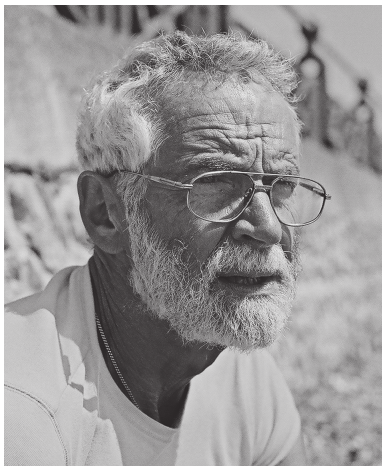
The algorithm is called the QR algorithm, invented independently by Francis (1959) and Kublanovskaya (1961). It is widely regarded as one of the ten most important algorithms of the 20^{th} century.



John Francis, 1934–

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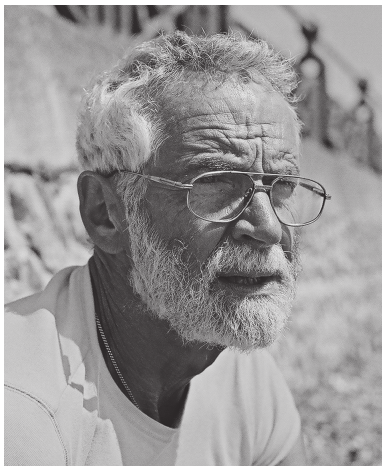

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You can learn more in A7: Numerical Analysis.

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be our (monic) polynomial. Then we can construct its *companion matrix*

$$C(a) := \begin{pmatrix} 0 & & & & -a_0 \\ 1 & 0 & & & -a_1 \\ 0 & 1 & 0 & & -a_2 \\ & & \ddots & & \vdots \\ & & & 1 & -a_{n-1} \end{pmatrix}.$$

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By construction, we have (proof is by induction):

$$\det(C(a) - \lambda I) = (-1)^n p(\lambda).$$

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We previously saw that Newton's method can get stuck in a cycle for $p(x) = x^3 - 2x + 2$. No problem:

```
>> roots([1 0 -2 2])

ans =

 -1.769292354238631 + 0.000000000000000i
  0.884646177119316 + 0.589742805022205i
  0.884646177119316 - 0.589742805022205i
```

Section 8

Bonus: representing polynomials

Philosophical remark

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For example, as mathematicians we might think of $p \in \Pi_n$, the vector space of degree-n polynomials. But Horner's method and the companion matrix *rely* on a particular representation of p, in the monomial basis $\{M_i\}$:

$$p(x) = \sum_{i=0}^{n} a_i M_i(x), \quad M_i(x) := x^i.$$

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A natural question to ask:

is the map $a \mapsto p$ *stable*?

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For the monomial basis $\{M_i\}$, the answer is *very very big*:

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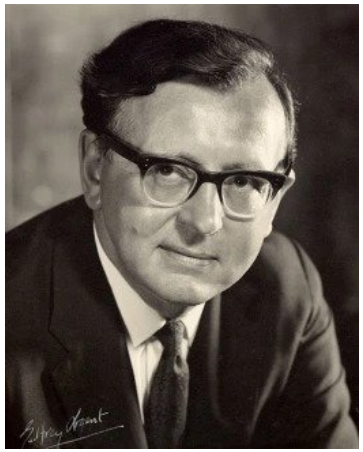
Construct

$$p(x) = \prod_{i=1}^{20}(x - i), \quad x \in [0, 20],$$

then perturb its monomial coefficients by

$$\delta a = [0, -2^{-23}, 0, \ldots, 0].$$



James H. Wilkinson, 1919–1986

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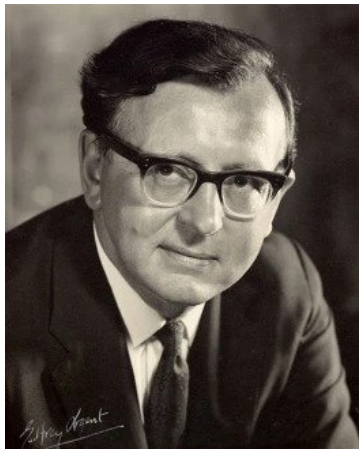

James H. Wilkinson, 1919–1986

The resulting δp has

$$\|\delta p\|_\infty := \max\{|\delta p(x)| : x \in [0, 20]\} \approx 6.25 \times 10^{17}$$

for a *stability constant* of

$$\frac{\|\delta p\|_\infty}{\|\delta a\|_\infty} \approx 5 \times 10^{24}.$$

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But you'd much rather compute with the latter than the former for small ε.

So what is a good basis for polynomials? An excellent choice on $[a, b]$ is

$$p(x) = \sum_{i=0}^{n} c_i T_i(\hat{x}(x)), \quad \hat{x} = \frac{2(x-a)}{(b-a)} - 1$$

where the *Chebyshev polynomials* $\{T_i : [-1, 1] \to [-1, 1]\}$ satisfy

$$T_0(\hat{x}) = 1, \quad T_1(\hat{x}) = \hat{x}, \quad T_{i+1}(\hat{x}) = 2\hat{x}T_i(\hat{x}) - T_{i-1}(\hat{x}).$$

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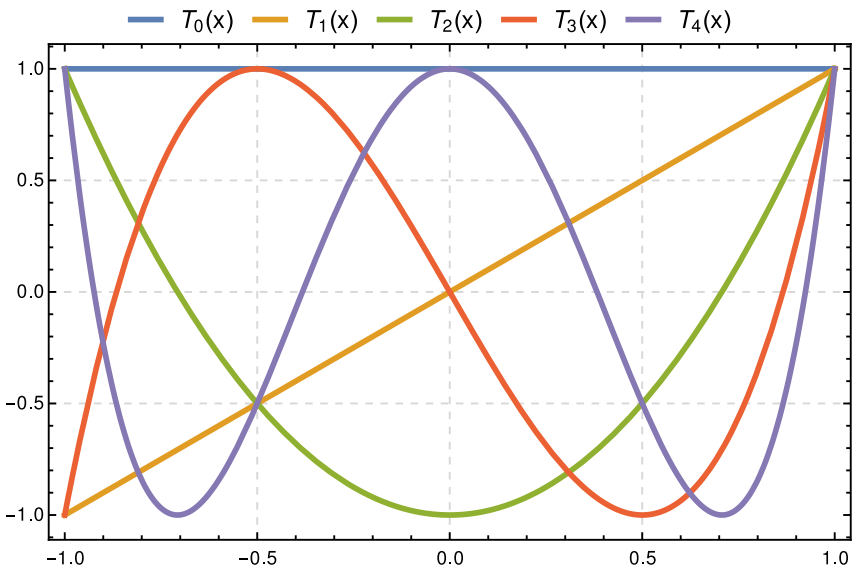

Chebyshev polynomials. Credit: Glosser.ca, Wikipedia

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Theorem

Let $f : [a, b] \to \mathbb{R}$ be analytic with Chebyshev expansion

$$f(x) = \sum_{i=0}^{\infty} c_i T_i(x).$$

Then for a constant $C > 1$

$$\|f - p_n\|_\infty = \mathcal{O}(C^{-n}), \quad p_n(x) = \sum_{i=0}^{n} c_i T_i(x).$$

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You can learn more in C6.3 Approximation of Functions.

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Good news

For Chebyshev bases, analogous algorithms exist:

- ✓ the *second barycentric formula*, for $\mathcal{O}(n)$ evaluation, and
- ✓ the *colleague matrix*, for finding all roots with the QR algorithm.

These allow us to work with polynomials with degrees in the millions.

M4: Constructive Mathematics
Lecture 4: Higher-dimensional rootfinding

Patrick E. Farrell

University of Oxford

We have considered several algorithms for rootfinding over $\mathbb{R}$:

$$\text{given } f : \mathbb{R} \to \mathbb{R}, \text{ find } x^\star \in \mathbb{R} \text{ such that } f(x^\star) = 0.$$

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- bisection ($q = 1$, $\mu = 1/2$, when it applies)
- secant method ($q = \phi \approx 1.618$, usually)
- Newton's method ($q = 2$, usually)
- Halley's method ($q = 3$, usually)

In real life, most problems involve more than one variable. So let's consider

$$\text{given } F : \mathbb{R}^N \to \mathbb{R}^N, \text{ find } \mathbf{x}^\star \in \mathbb{R}^N \text{ such that } F(\mathbf{x}^\star) = \mathbf{0}.$$

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Simpson extended Newton's method to this case in his 1740 book *Essays on Several Curious and Useful Subjects in Speculative and Mix'd Mathematicks, Illustrated by a Variety of Examples*.

Thomas Simpson, 1710–1761

Section 2

Derivation of Newton's method

The geometric pictures we had in one dimension don't naturally extend to higher dimensions. So first let's see another derivation of Newton's method in $\mathbb{R}$ that does extend.

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Consider a Taylor expansion of f. We want to find $x_{i+1} = x_i + \delta x$:

$$f(x_i + \delta x) = f(x_i) + \delta x f'(x_i) + \text{higher-order terms.}$$

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We want to choose the update δx so that $f(x_i + \delta x) = 0$. Setting the left-hand side to zero, and dropping higher-order terms, we get

$$\delta x = -[f'(x_i)]^{-1} f(x_i), \quad x_{i+1} = x_i + \delta x,$$

which we recognise as Newton's scheme written in update form.

Taylor's theorem extends to higher dimensions, with the role of derivative f' replaced by the **Jacobian** matrix. If $F : \mathbb{R}^N \to \mathbb{R}^N$ looks like

$$F(\mathbf{x}) = F\begin{pmatrix} \mathbf{x}^1 \\ \mathbf{x}^2 \\ \vdots \\ \mathbf{x}^N \end{pmatrix} = \begin{pmatrix} F^1(\mathbf{x}^1, \ldots, \mathbf{x}^N) \\ F^2(\mathbf{x}^1, \ldots, \mathbf{x}^N) \\ \vdots \\ F^N(\mathbf{x}^1, \ldots, \mathbf{x}^N) \end{pmatrix},$$

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then its Jacobian $DF : \mathbb{R}^N \to \mathbb{R}^{N \times N}$ is

$$DF(\mathbf{a}) := \begin{bmatrix} \frac{\partial F^1}{x^1}(\mathbf{a}) & \frac{\partial F^1}{x^2}(\mathbf{a}) & \cdots & \frac{\partial F^1}{x^N}(\mathbf{a}) \\ \frac{\partial F^2}{x^1}(\mathbf{a}) & \frac{\partial F^2}{x^2}(\mathbf{a}) & \cdots & \frac{\partial F^2}{x^N}(\mathbf{a}) \\ \vdots & \vdots & & \vdots \\ \frac{\partial F^N}{x^1}(\mathbf{a}) & \frac{\partial F^N}{x^2}(\mathbf{a}) & \cdots & \frac{\partial F^N}{x^N}(\mathbf{a}) \end{bmatrix}.$$

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Following the reasoning from one dimension,

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In practice, we don't actually invert the matrix, but rather

$$\text{solve } DF(\mathbf{x}_i)\delta\mathbf{x} = -F(\mathbf{x}_i),$$

using e.g. an LU factorisation of the matrix.

Newton–Raphson method

$$\mathbf{x}_{i+1} = g(\mathbf{x}_i) := \mathbf{x}_i - (DF(\mathbf{x}_i))^{-1} F(\mathbf{x}_i).$$

Comments:

- ✓ Still a fixed-point method.

Newton–Raphson method

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- ✗ If x_0 is far away, the method can diverge or get stuck in a cycle.
- ✓ Newton's method even generalises to infinite dimensions.

Section 3

Example

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Section 4

Convergence

Definition (Norm of $\mathbf{x} \in \mathbb{R}^N$)

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Definition (Order of convergence of a sequence)

Suppose $(\mathbf{x}_i) \to \mathbf{x}^\star$. The sequence converges with order q if

$$\lim_{i\to\infty} \frac{\|\mathbf{x}_{i+1} - \mathbf{x}^\star\|_\infty}{\|\mathbf{x}_i - \mathbf{x}^\star\|_\infty^q} = M$$

for some $M > 0$ (if $q = 1$ we need $M < 1$).

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Theorem (Quadratic convergence of Newton's method)

Let $F \in C^2(\mathbb{R}^N, \mathbb{R}^N)$, i.e. F is continuous with all first and second partial derivatives continuous. Suppose $\mathbf{x}^\star \in \mathbb{R}^N$ is an isolated root of F, i.e. $F(\mathbf{x}^\star) = \mathbf{0}$ with $DF(\mathbf{x}^\star)$ nonsingular. Then if $\mathbf{x}_0$ is close enough to $\mathbf{x}^\star$, the Newton sequence will converge quadratically.

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The core of the proof is that the Jacobian matrix of the associated fixed-point iteration is zero at $\mathbf{x}^\star$.

Section 5

Bonus: Affine covariance

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Now imagine that we change units or coordinate systems for our outputs F. Instead of solving $F(\mathbf{x}) = \mathbf{0}$, we want to solve $\tilde{F}(\mathbf{x}) = AF(\mathbf{x}) = \mathbf{0}$, where $A \in \mathbb{R}^{N \times N}$ is constant and nonsingular. Of course, this doesn't change the roots $\mathbf{x}^\star$.

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Theorem (Affine covariance)

Premultiplying F by a constant nonsingular $A \in \mathbb{R}^{N \times N}$ does not change the Newton sequence.

Let $\tilde{F}(\mathbf{x}) := AF(\mathbf{x})$. Newton's method applied to $\tilde{F}$ from $\mathbf{x}_0 = \tilde{\mathbf{x}}_0$ generates a sequence

$$\tilde{\mathbf{x}}_0, \tilde{\mathbf{x}}_1, \tilde{\mathbf{x}}_2, \ldots.$$

Proof.

For $i = 0$, we have $\mathbf{x}_i = \tilde{\mathbf{x}}_i$ by assumption.

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We get exactly the same iterates $\mathbf{x}_0, \mathbf{x}_1, \ldots$, whether we apply Newton to $F(\mathbf{x}) = \mathbf{0}$ or $AF(\mathbf{x}) = \mathbf{0}$.

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Moreover, any sensible strategy for globalising the convergence of Newton's method from poor initial guesses $\mathbf{x}_0$ must also preserve this property. This insight leads to the current state of the art for globalising Newton's method.

Peter Deuflhard, 1944–2019

Section 6

Bonus: the Newton–Kantorovich theorem

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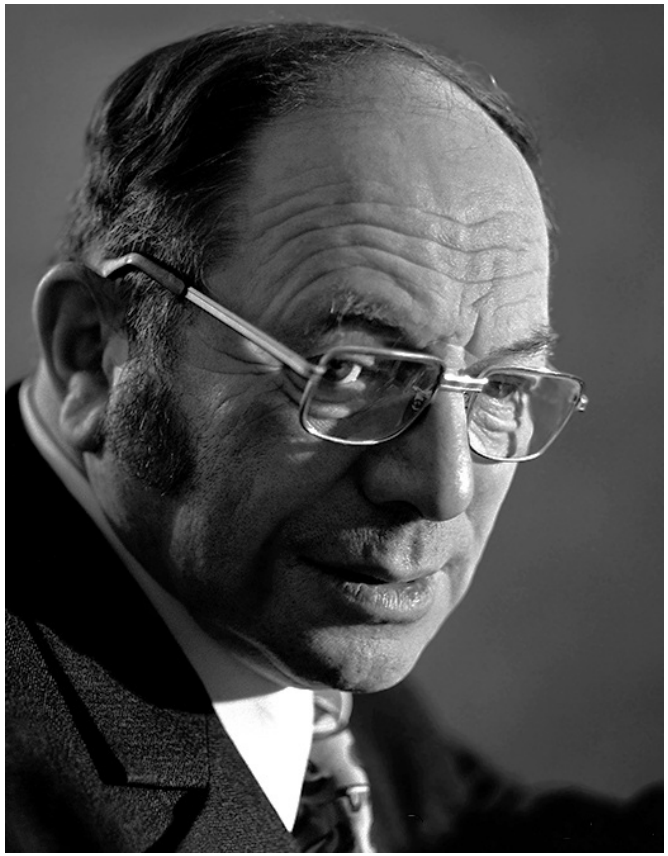

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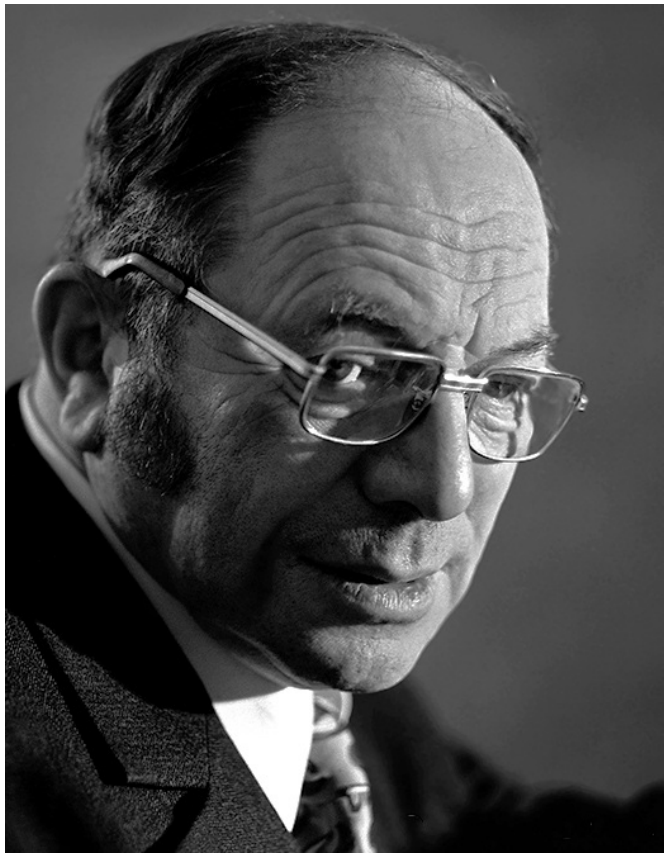

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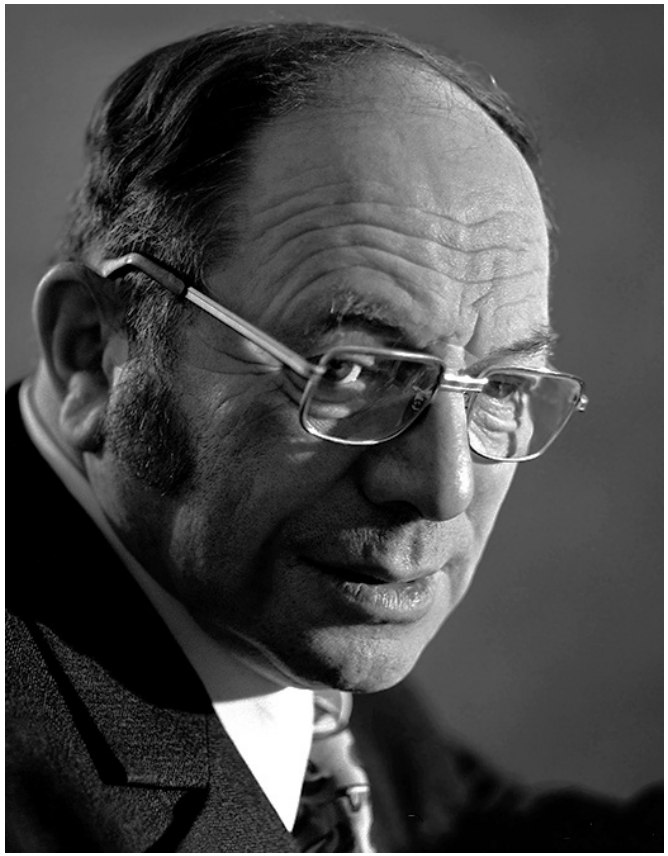

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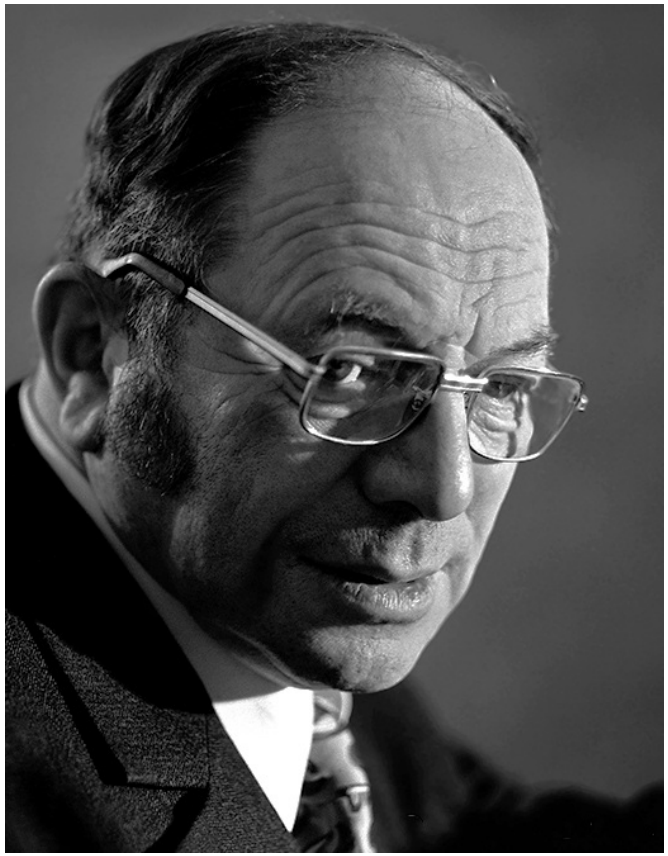

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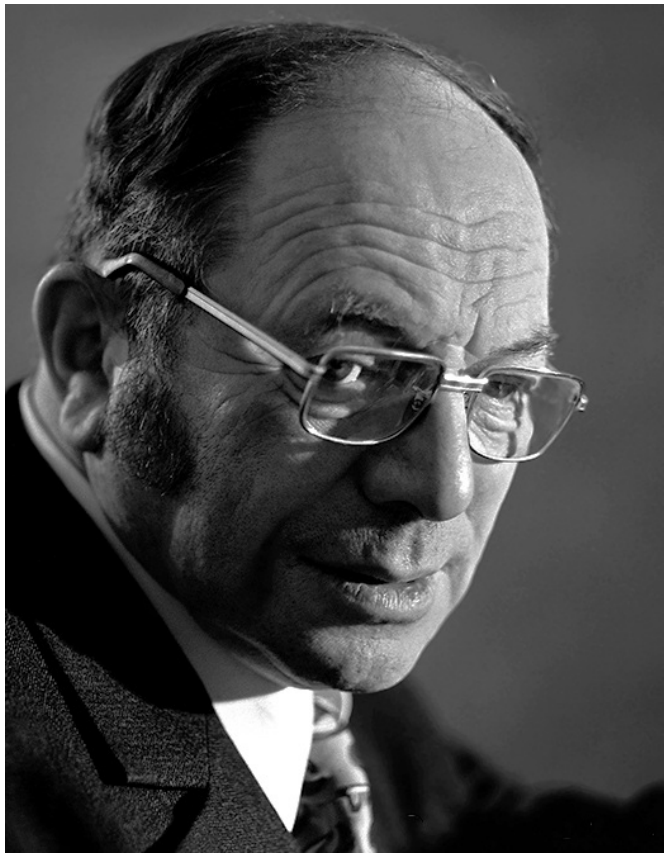

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With a good initial guess, and great cleverness, it is possible to devise *computer-assisted proofs* of the existence of solutions to infinite-dimensional nonlinear problems.

Theorem (Kantorovich (1948))

Let $F \in C^1(\mathbb{R}^N, \mathbb{R}^N)$ be the residual of our nonlinear problem, and let $\mathbf{x}_0 \in \mathbb{R}^N$ be an initial guess such that the Fréchet derivative (Jacobian) $DF(\mathbf{x}_0)$ is invertible. Let $B(\mathbf{x}_0, r)$ denote the open ball of radius r centred at $\mathbf{x}_0$.

Assume that there exists a constant $r > 0$ such that

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(2) For all $\tilde{\mathbf{x}}, \mathbf{x} \in B(\mathbf{x}_0, r)$,

$$\|DF(\mathbf{x}_0)^{-1} \left(DF(\tilde{\mathbf{x}}) - DF(\mathbf{x})\right)\| \leq \frac{1}{r} \|\tilde{\mathbf{x}} - \mathbf{x}\|.$$

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(3) For each $i \geq 0$,

$$\|\mathbf{x}^\star - \mathbf{x}_i\| \leq \frac{r}{2^i}.$$

Theorem (Kantorovich (1948))

Then

(1) $DF(\mathbf{x})$ *is invertible at each* $\mathbf{x} \in B(\mathbf{x}_0, r)$.

(2) The Newton sequence $(\mathbf{x}_i)_{i=0}^{\infty}$ *defined by*

$$\mathbf{x}_{i+1} = \mathbf{x}_i - DF(\mathbf{x}_i)^{-1} F(\mathbf{x}_i)$$

satisfies $\mathbf{x}_i \in B(\mathbf{x}_0, r)$ *for all* i*, and converges to a root* $\mathbf{x}^\star$ *of* F.

(3) For each $i \geq 0$,

$$\|\mathbf{x}^\star - \mathbf{x}_i\| \leq \frac{r}{2^i}.$$

(4) The root $\mathbf{x}^\star$ *is locally unique, i.e.* $\mathbf{x}^\star$ *is the only root of* F *in the ball* $\overline{B(\mathbf{x}_0, r)}$.

Section 7

Bonus: The Davidenko differential equation

Newton's method applied to $F(\mathbf{x}) = \mathbf{0}$ produces a sequence

$$\mathbf{x}_0, \mathbf{x}_1, \mathbf{x}_2, \ldots, \quad \mathbf{x}_i \in \mathbb{R}^N.$$

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Yes. The *Davidenko differential equation* is

$$\frac{\mathrm{d}\mathbf{x}}{\mathrm{d}s} = -[DF(\mathbf{x})]^{-1} F(\mathbf{x}).$$

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The Newton iteration is the forward Euler discretisation of the Davidenko differential equation with $\Delta s = 1$:

$$\frac{\mathrm{d}\mathbf{x}}{\mathrm{d}s} \approx \frac{\mathbf{x}(s + \Delta s) - \mathbf{x}(s)}{\Delta s} = -[DF(\mathbf{x}(s))]^{-1} F(\mathbf{x}(s)).$$

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This shows us that the tangent of the curve—the Newton update $[DF(\mathbf{x})]^{-1}F(\mathbf{x})$—is a special direction to go to find a root, even far away from a solution. It's just that it might be *too long*.

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You can use these ideas to build effective globalisation strategies for Newton's method.

Section 8

Newton fractals

One last beautiful idea about Newton's method in higher dimensions.

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Consider the problem

$$\text{find } z \in \mathbb{C} \text{ such that } z^3 - 1 = 0.$$

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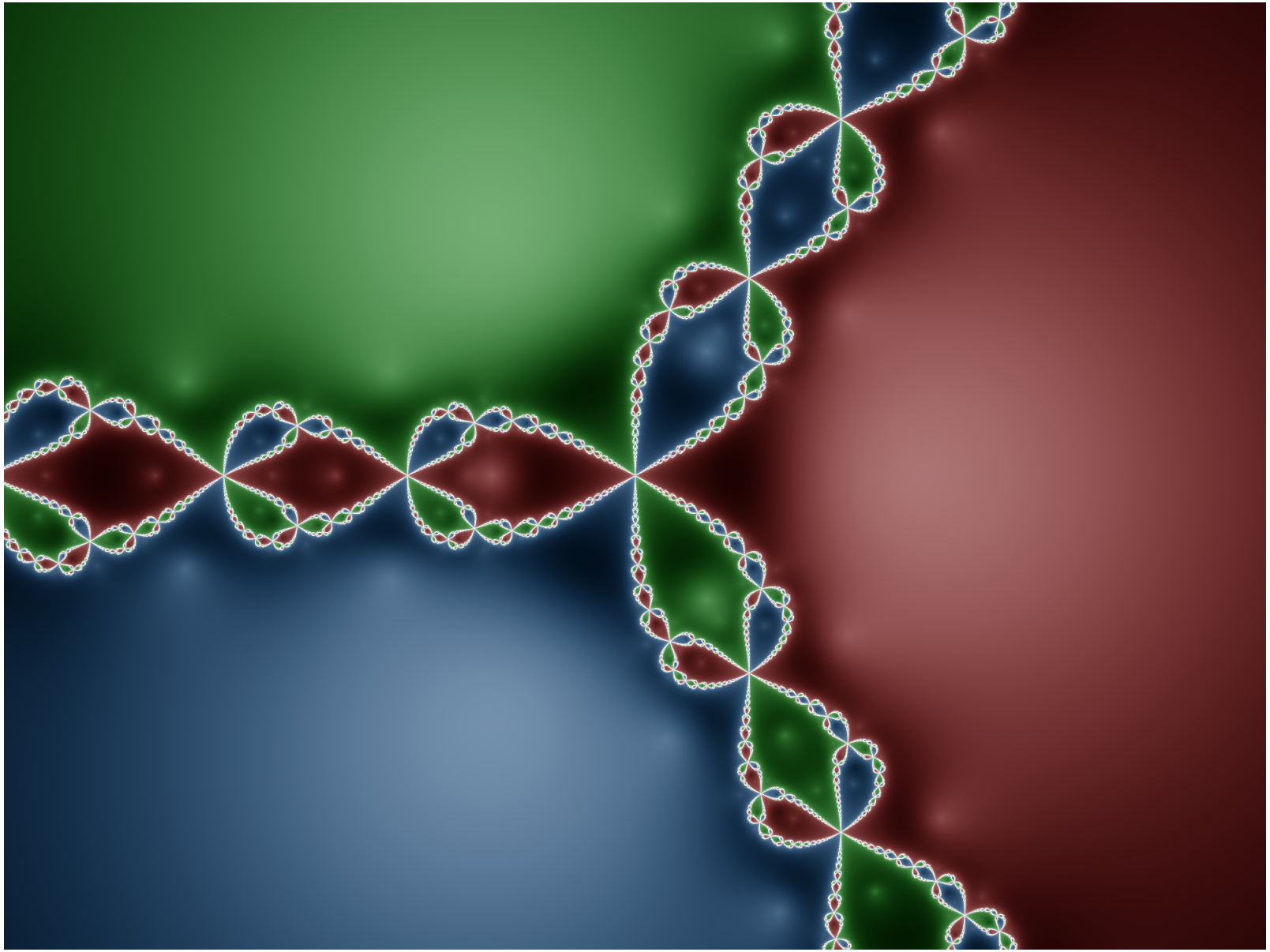
We could also think of this as a problem in $\mathbb{R}^2$.

We know this has three solutions,

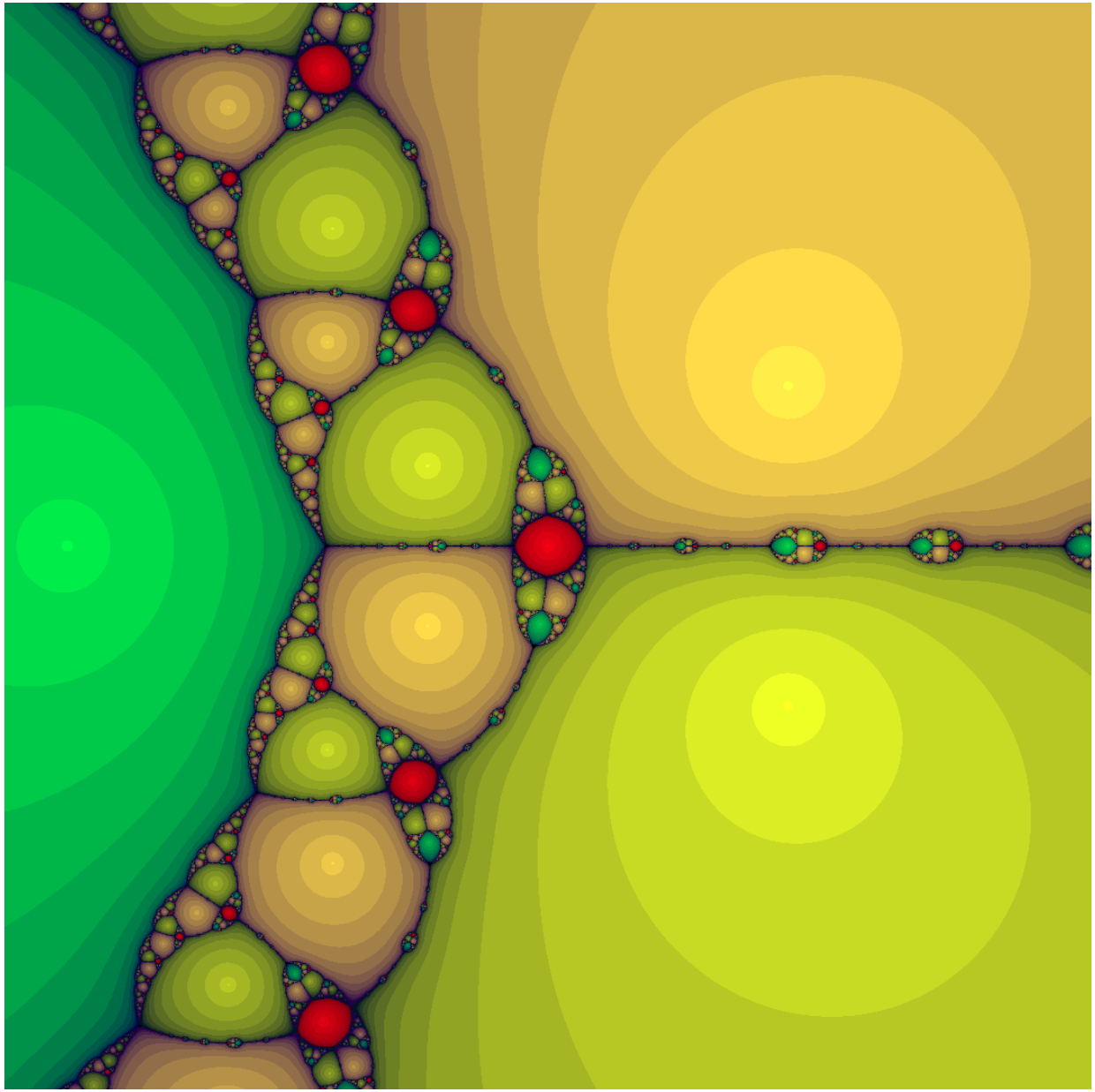
$$z = 1, \quad z = -1/2 + i\sqrt{3}/2, \text{ and } z = -1/2 - i\sqrt{3}/2.$$

Let's take a subset of the complex plane and colour each point as follows. For a given $z_0 \in \mathbb{C}$, we

1. run Newton's method with that initial guess,
2. colour the point according to which root it converges to,
3. shade the colour by how many iterations it took.



The Newton fractal for $z^3 - 1 = 0$.



The Newton fractal for $z^3 - 2z + 2 = 0$.

Section 9

Algorithms for optimisation problems

In this final lecture, we study how to apply rootfinding ideas to *optimisation*.

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The ideas in this lecture are further explored in ASO Calculus of Variations, B6.2 Optimisation for Data Science (new!), and C6.2 Continuous Optimisation.

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Let's consider the optimisation problem: given $f \in C^2(\mathbb{R}^N, \mathbb{R})$,

$$\text{find } \mathbf{x}^\star = \operatorname*{argmin}_{\mathbf{x} \in \mathbb{R}^N} f(\mathbf{x}).$$

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We want an argument $\mathbf{x}^\star$ that satisfies $f(\mathbf{x}^\star) \leq f(\mathbf{x})$ for all $\mathbf{x} \in \mathbb{R}^N$.

This is usually too much to ask for, so instead we satisfy ourselves with *local minima* $\mathbf{x}^\star$ such that there is a neighbourhood $\mathcal{N}$ around $\mathbf{x}^\star$ so that

$$f(\mathbf{x}^\star) \leq f(\mathbf{x}) \text{ for all } \mathbf{x} \in \mathcal{N}.$$

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...and local maxima, and saddle points: any point satisfying $g(\mathbf{x}) = \mathbf{0}$ is called a *critical point*.

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Finding global minimisers is so hard that it is its own branch of study, *global* optimisation.

The model problem we're considering in this lecture is quite simplified. In most real optimisation problems, there are *constraints* on the solution:

$$\begin{aligned} \min_{\mathbf{x}\in\mathbb{R}^N} \quad & f(\mathbf{x}) \\ \text{subject to} \quad & c_i(\mathbf{x}) \geq 0, \quad i \in \mathcal{I}, \\ & c_e(\mathbf{x}) = 0, \quad i \in \mathcal{E}. \end{aligned}$$

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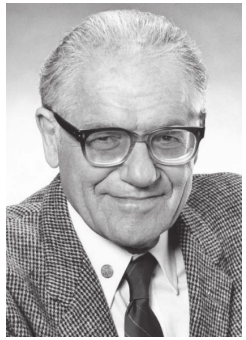
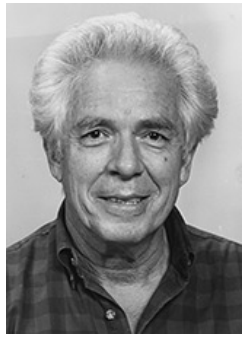
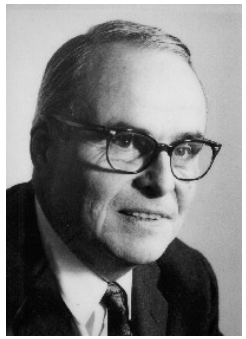
For problems with constraints, the optimality conditions are no longer as simple as $\nabla f(\mathbf{x}) = 0$. The optimality conditions for the problem above are known as the *Karush–Kuhn–Tucker* conditions.



William Karush, 1917–1997



Harold Kuhn, 1925–2014



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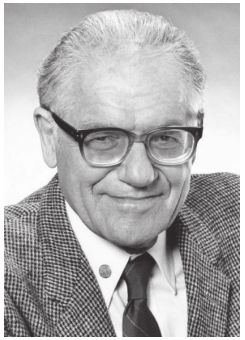
In this lecture we consider the unconstrained problem, since you need to understand that first to attack the constrained one!



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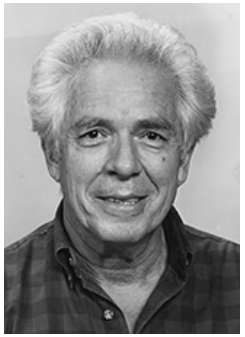

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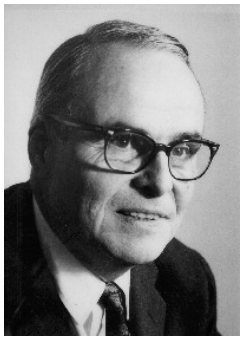

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Section 10

Newton's method for optimisation

Let's see what Newton iteration on the gradient looks like. If we take the Jacobian of the gradient, we get the *Hessian matrix*:

$$Hf(\mathbf{a}) = D\nabla f(\mathbf{a}) := \begin{pmatrix} \frac{\partial^2 f}{\partial x^1 x^1}(\mathbf{a}) & \frac{\partial^2 f}{\partial x^1 x^2}(\mathbf{a}) & \cdots & \frac{\partial^2 f}{\partial x^1 x^N}(\mathbf{a}) \\ \frac{\partial^2 f}{\partial x^2 x^1}(\mathbf{a}) & \frac{\partial^2 f}{\partial x^2 x^2}(\mathbf{a}) & \cdots & \frac{\partial^2 f}{\partial x^2 x^N}(\mathbf{a}) \\ \vdots & \vdots & & \vdots \\ \frac{\partial^2 f}{\partial x^N x^1}(\mathbf{a}) & \frac{\partial^2 f}{\partial x^N x^2}(\mathbf{a}) & \cdots & \frac{\partial^2 f}{\partial x^N x^N}(\mathbf{a}) \end{pmatrix}.$$

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$$\mathbf{x}_{i+1} = \mathbf{x}_i - [Hf(\mathbf{x}_i)]^{-1} \nabla f(\mathbf{x}_i) = x_i - [Dg(\mathbf{x}_i)]^{-1} g(\mathbf{x}_i).$$

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Suppose we're at iterate $\mathbf{x}_i$ and we'd like to minimise f. We don't know how, so we'll replace f with a local quadratic model:

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So at every step, Newton's method for optimisation approximates the function with a paraboloid, and minimises that.

Section 11

Quasi-Newton methods

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2. How do we store $Hf(\mathbf{x}_i)$? (Can't store a full/dense matrix.)
3. How do we solve $Hf(\mathbf{x}_i)\delta\mathbf{x} = -\nabla f(\mathbf{x}_i)$?

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It is often possible to overcome these issues by exploiting some *structure* in the problem. When minimising energy functions in physics, the matrix is usually *sparse*, which can sometimes be exploited to solve the linear system in time $\mathcal{O}(N)$ instead of $\mathcal{O}(N^3)$.

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The standard practice is to modify the algorithm to

$$\mathbf{x}_{i+1} = \mathbf{x}_i - B_i^{-1}\nabla f(\mathbf{x}_i)$$

for carefully chosen matrices B_i. This is called a *quasi-Newton* scheme.

Here are some choices for B_i:

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2. $B_i = I \in \mathbb{R}^{N \times N}$. Gradient descent. Linear convergence, very slow.

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The choice most used in practice is the BFGS algorithm (1970).

This builds up an approximation to the Hessian as the iterations proceed.

The BFGS approach demands that the symmetric matrix B_{i+1} satisfy

$$B_{i+1}(\mathbf{x}_{i+1} - \mathbf{x}_i) = \nabla f(\mathbf{x}_{i+1}) - \nabla f(\mathbf{x}_i).$$

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This means we now need to supply B_0. With the right choice of norm, this problem has an explicit solution for B_{i+1} and B_{i+1}^{-1}.

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Diagonal matrices

A diagonal matrix A is positive-definite iff all of its diagonal entries are strictly positive. In this case,

$$\mathbf{x}^T A\mathbf{x} = A_{11}(\mathbf{x}^1)^2 + A_{22}(\mathbf{x}^2)^2 + \cdots + A_{NN}(\mathbf{x}^N)^2 > 0.$$

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BFGS gives a positive-definite Hessian approximation, if B_0 is.

To ensure we satisfy

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and consider its derivative at $t = 0$:

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We therefore modify the algorithm to

$$\mathbf{x}_{i+1} = \mathbf{x}_i - t_i^\star B_i^{-1} \nabla f(\mathbf{x}_i),$$

where $t_i^\star$ is an (approximate) minimiser of $\phi(t)$.

We end the course with a final example. Consider the problem

$$\text{find } (x,y)^\star = \operatorname*{argmin}_{(x,y)\in\mathbb{R}^2} f(x,y) := 100(y-x^2)^2 + (1-x)^2.$$

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Gradient descent took 5264 iterations, Newton's method 21, and BFGS 34.