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

Patrick E. Farrell

University of Oxford

For $a_0, a_1 \in \mathbb{R}, a_1 \neq 0$, find $x \in \mathbb{R}$ such that $a_1x + a_0 = 0$.

For $a_0, a_1 \in \mathbb{R}, a_1 \neq 0$, find $x \in \mathbb{R}$ such that $a_1x + a_0 = 0$.

This has solution $x = -a_0/a_1$.

For $a_0, a_1 \in \mathbb{R}, a_1 \neq 0$, find $x \in \mathbb{R}$ such that $a_1x + a_0 = 0$.

This has solution $x = -a_0/a_1$.

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

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

For $a_0, a_1 \in \mathbb{R}, a_1 \neq 0$, find $x \in \mathbb{R}$ such that $a_1x + a_0 = 0$.

This has solution $x = -a_0/a_1$.

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

For $a_0, \ldots, a_5 \in \mathbb{R}$, find $x \in \mathbb{C}$ such that $a_5x^5 + a_4x^4 + \cdots + a_0 = 0$.

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

So what do we do in this situation? We still care about the roots of polynomials!

So what do we do in this situation? We still care about the roots of polynomials!

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 $(-1)^n a_0/a_n$.

So what do we do in this situation? We still care about the roots of polynomials!

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 $(-1)^n a_0/a_n$.

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

Algorithms are of core interest to both pure and applied mathematics.

Algorithms are of core interest to both pure and applied mathematics.

In pure mathematics, we use algorithms to (among other things) prove the existence of some object. We will see examples in the course.

Algorithms are of core interest to both pure and applied mathematics.

In pure mathematics, we use algorithms to (among other things) prove the existence of some object. We will see examples in the course.

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

find
$$y(t)$$
 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 $y_n \rightarrow y$.

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



P. E. Farrell (Oxford)

What is constructive mathematics?

Does our algorithm terminate?

Does our algorithm terminate?

Theorem (Halting problem, 1936)

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



Alan Turing, 1912-1954

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

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

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

How fast does the algorithm converge to the right answer?

How fast does the algorithm converge to the right answer?

Consider two formulae for π :

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

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



Gottfried Leibniz, 1646–1716



Srinivasa Ramanujan, 1887–1920

How fast does the algorithm converge to the right answer?

Consider two formulae for π :

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

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

About 5 billion, vs 2!



Gottfried Leibniz, 1646–1716



Srinivasa Ramanujan, 1887–1920

How many operations does the algorithm take?

How many operations does the algorithm take?

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

Are the problem and algorithm stable to small perturbations in data?

Are the problem and algorithm stable to small perturbations in data?

Consider $p(x) = (x - 1)(x - 2) \cdots (x - 20)$. Expanding in monomials, we have

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



James H. Wilkinson, 1919-1986

Are the problem and algorithm stable to small perturbations in data?

Consider $p(x) = (x - 1)(x - 2) \cdots (x - 20)$. Expanding in monomials, we have

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

If we perturb -210 to -210.0000001192, then the roots become (to 5 digits)



James H. Wilkinson, 1919-1986

1.00000 2.00000 3.00000 4.00000 5.00000

Are the problem and algorithm stable to small perturbations in data?

Consider $p(x) = (x - 1)(x - 2) \cdots (x - 20)$. Expanding in monomials, we have

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

If we perturb -210 to -210.0000001192, then the roots become (to 5 digits)



James H. Wilkinson, 1919-1986

| 1.00000 | 2.00000 | 3.00000 | 4.00000 | 5.00000 |
|---------|---------|---------|---------|----------|
| 6.00001 | 6.99970 | 8.00727 | 8.91725 | 20.84691 |

Are the problem and algorithm stable to small perturbations in data?

Consider $p(x) = (x - 1)(x - 2) \cdots (x - 20)$. Expanding in monomials, we have

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

If we perturb -210 to -210.0000001192, then the roots become (to 5 digits)



James H. Wilkinson, 1919-1986

| 1.00000 | 2.00000 | 3.00000 | 4.00000 | 5.00000 |
|--|-------------------------------|--------------------------|---------------------------|--|
| 6.00001 | 6.99970 | 8.00727 | 8.91725 | 20.84691 |
| $\begin{array}{c} 10.09527 \pm \\ 0.64350 i \end{array}$ | ${}^{11.79363\pm}_{1.65233i}$ | $13.99236 \pm 2.51883 i$ | ${16.73074\pm}{2.81262i}$ | $\begin{array}{c} 19.50244 \pm \\ 1.94033 i \end{array}$ |

M4: Constructive Mathematics Lecture 1: Euclid's algorithm

Patrick E. Farrell

University of Oxford

We start with the natural numbers

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

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

$$t = qb + r, \quad 0 \le r < b.$$

We start with the natural numbers

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

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

$$t = qb + r, \quad 0 \le r < b.$$

Either b divides t evenly (r = 0), or we are allowed to divide b by r. If $r \neq 0$ we can *iterate*: we can see what the division of b by r is.

We start with the natural numbers

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

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

$$t = qb + r, \quad 0 \le r < b.$$

Either b divides t evenly (r = 0), or we are allowed to divide b by r. If $r \neq 0$ we can *iterate*: we can see what the division of b by r is.

Take t = 80, b = 7:

$$80 = 11 \times 7 + 3 \quad (q = 11, r = 3)$$

We start with the natural numbers

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

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

$$t = qb + r, \quad 0 \le r < b.$$

Either b divides t evenly (r = 0), or we are allowed to divide b by r. If $r \neq 0$ we can *iterate*: we can see what the division of b by r is.

Take t = 80, b = 7: $80 = 11 \times 7 + 3 \quad (q = 11, r = 3)$ $7 = 2 \times 3 + 1 \quad (q = 2, r = 1)$

We start with the natural numbers

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

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

$$t = qb + r, \quad 0 \le r < b.$$

Either b divides t evenly (r = 0), or we are allowed to divide b by r. If $r \neq 0$ we can *iterate*: we can see what the division of b by r is.

Take t = 80, b = 7: $80 = 11 \times 7 + 3 \quad (q = 11, r = 3)$ $7 = 2 \times 3 + 1 \quad (q = 2, r = 1)$

We start with the natural numbers

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

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

$$t = qb + r, \quad 0 \le r < b.$$

Either b divides t evenly (r = 0), or we are allowed to divide b by r. If $r \neq 0$ we can *iterate*: we can see what the division of b by r is.

Take t = 80, b = 7:

$$\begin{array}{ll} 80 = 11 \times 7 + 3 & (q = 11, r = 3) \\ 7 = & 2 \times 3 + 1 & (q = & 2, r = 1) \\ 3 = & 3 \times 1 + 0 & (q = & 3, r = 0). \end{array}$$

We start with the natural numbers

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

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

$$t = qb + r, \quad 0 \le r < b.$$

Either b divides t evenly (r = 0), or we are allowed to divide b by r. If $r \neq 0$ we can *iterate*: we can see what the division of b by r is.

Take t = 80, b = 7:

$$\begin{array}{ll} 80 = 11 \times 7 + 3 & (q = 11, r = 3) \\ 7 = & 2 \times 3 + 1 & (q = & 2, r = 1) \\ 3 = & 3 \times 1 + 0 & (q = & 3, r = 0). \end{array}$$

The game ends when r = 0.

We start with the natural numbers

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

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

$$t = qb + r, \quad 0 \le r < b.$$

Either b divides t evenly (r = 0), or we are allowed to divide b by r. If $r \neq 0$ we can *iterate*: we can see what the division of b by r is.

Take t = 80, b = 7:

$$80 = 11 \times 7 + 3 \quad (q = 11, r = 3)$$

$$7 = 2 \times 3 + 1 \quad (q = 2, r = 1)$$

$$3 = 3 \times 1 + 0 \quad (q = 3, r = 0).$$

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 \mod b

while r \neq 0 do

t \leftarrow b

b \leftarrow r

r \leftarrow t \mod b

end while

return b

end function
```

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 \mod b

while r \neq 0 do

t \leftarrow b

b \leftarrow r

r \leftarrow t \mod b

end while

return b

end function
```

This completely and unambiguously lists the steps for a computer to take.

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 \mod b

while r \neq 0 do

t \leftarrow b

b \leftarrow r

r \leftarrow t \mod b

end while

return b

end function
```

This completely and unambiguously lists the steps for a computer to take.

Note that this algorithm calls another one (the division algorithm).

P. E. Farrell (Oxford)

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

Given any $t, b \in \mathbb{N}$, 0 < b < t, Euclid's algorithm computes the greatest common divisor of t and b.



Euclid of Alexandria, c. 300 BCE

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

Given any $t, b \in \mathbb{N}$, 0 < b < t, Euclid's algorithm computes the greatest common divisor of t and b.

For convenience, let's label each intermediate value:

$$t = q_0 b + r_0$$

$$b = q_1 r_0 + r_1$$

$$r_0 = q_2 r_1 + r_2$$

:

$$r_j = q_{j+2} r_{j+1} + r_{j+2}$$

:



Euclid of Alexandria, c. 300 BCE

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

Given any $t, b \in \mathbb{N}$, 0 < b < t, Euclid's algorithm computes the greatest common divisor of t and b.

For convenience, let's label each intermediate value:

$$t = q_0 b + r_0$$

$$b = q_1 r_0 + r_1$$

$$r_0 = q_2 r_1 + r_2$$

:

$$r_j = q_{j+2} r_{j+1} + r_{j+2}$$

:



Euclid of Alexandria, c. 300 BCE

Also for convenience, denote

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

Claim: the algorithm terminates.

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

Claim: the algorithm terminates.

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

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.

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

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.

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

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

$$r_{i-3} = q_{i-1}r_{i-2} + r_{i-1}$$
$$= (\cdots) \times r_{i-1}$$

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.

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

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

$$r_{i-3} = q_{i-1}r_{i-2} + r_{i-1} \\ = (\cdots) \times r_{i-1}$$

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.

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

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

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

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

Since $t = q_0 b + r_0$, we get $r_0 = (\alpha - q_0 \beta)d$, so d divides r_0 .

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

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

Since $t = q_0 b + r_0$, we get $r_0 = (\alpha - q_0 \beta)d$, so d divides r_0 .

The next equation is $b = q_1r_0 + r_1$, 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} .

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

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

Since $t = q_0 b + r_0$, we get $r_0 = (\alpha - q_0 \beta)d$, so d divides r_0 .

The next equation is $b = q_1r_0 + r_1$, 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} .

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

In this context, we ask: can we bound the number of divisions required in computing gcd(t, b) in terms of t and b, t > b > 0?

In this context, we ask: can we bound the number of divisions required in computing gcd(t, b) in terms of t and b, t > b > 0?

Since the remainder decreases at each iteration, we know at least that we will do at most b iterations, i.e. the cost grows linearly in the size of the inputs.

In this context, we ask: can we bound the number of divisions required in computing gcd(t, b) in terms of t and b, t > b > 0?

Since the remainder decreases at each iteration, we know at least that we will do at most b iterations, i.e. the cost grows linearly in the size of the inputs.

But it is possible to prove a tighter bound!

Theorem

Let t > b > 0. The smallest values of t and b for which Euclid's algorithm requires N iterations 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

Theorem

Let t > b > 0. The smallest values of t and b for which Euclid's algorithm requires N iterations 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.

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



Gabriel Lamé, 1795-1870



Greatest common divisors



Euclid's algorithm

Section 2

Diophantine equations

Some Diophantine equations have no solutions, like 4x + 6y = 3.

Some Diophantine equations have no solutions, like 4x + 6y = 3.

Some do, however. For example, 48x - 35y = 1 has a solution x = -8, y = -11.

Some Diophantine equations have no solutions, like 4x + 6y = 3.

Some do, however. For example, 48x - 35y = 1 has a solution x = -8, y = -11.

Diophantus' work was collected in his magnum opus, *Arithmetica*. In 1637, Pierre de Fermat wrote in the margin of his copy of *Arithmetica*,

It is impossible ... for any number which is a power greater than the second to be written as the sum of two like powers. I have a truly marvelous demonstration of this proposition which this margin is too narrow to contain.



Pierre de Fermat, 1607-1665

A linear Diophantine equation (LDE) in two variables is of the form: given $a, b, c \in \mathbb{Z}$, find $x, y \in \mathbb{Z}$ such that

$$ax + by = c.$$

A linear Diophantine equation (LDE) in two variables is of the form: given $a, b, c \in \mathbb{Z}$, find $x, y \in \mathbb{Z}$ such that

$$ax + by = c.$$

LDEs with gcd(a, b) = 1 = c are of particular interest. If we can solve

ax + by = 1

then we have solved the problem: find $x \in \mathbb{Z}$ such that

 $ax \equiv 1 \pmod{b},$

the problem of finding modular multiplicative inverses.

A linear Diophantine equation (LDE) in two variables is of the form: given $a, b, c \in \mathbb{Z}$, find $x, y \in \mathbb{Z}$ such that

$$ax + by = c.$$

LDEs with gcd(a, b) = 1 = c are of particular interest. If we can solve

ax + by = 1

then we have solved the problem: find $x \in \mathbb{Z}$ such that

 $ax \equiv 1 \pmod{b},$

the problem of finding modular multiplicative inverses.

In particular, this is a crucial step in RSA key generation: the private key \boldsymbol{d} satisfies

$$de \equiv 1 \pmod{\lambda(n)},$$

where n, e are the public key, and $\lambda(n)$ is easy to compute if you know the prime factorisation of n and difficult otherwise.

P. E. Farrell (Oxford)

Lemma (Bézout's Lemma)

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



Étienne Bézout, 1730-1783

Lemma (Bézout's Lemma)

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

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.



Étienne Bézout, 1730-1783



Claude Gaspar Bachet de Méziriac, 1581–1638

Lemma (Bézout's Lemma)

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

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.

Many other results in number theory follow from Bézout's Lemma, such as Euclid's Lemma and Sunzi's Remainder Theorem.



Étienne Bézout, 1730-1783



Claude Gaspar Bachet de Méziriac, 1581–1638

Before we prove Bézout's Lemma, let's do an example. Let's take 48x - 35y = 1 that we saw earlier. Applying Euclid's algorithm, we get

 $48 = 1 \times 35 + 13 \qquad 13 = 48 - 1 \times 35$

Before we prove Bézout's Lemma, let's do an example. Let's take 48x - 35y = 1 that we saw earlier. Applying Euclid's algorithm, we get

$$48 = 1 \times 35 + 13 \qquad 13 = 48 - 1 \times 35$$
$$35 = 2 \times 13 + 9 \qquad 9 = 35 - 2 \times 13$$
Before we prove Bézout's Lemma, let's do an example. Let's take 48x - 35y = 1 that we saw earlier. Applying Euclid's algorithm, we get

$$\begin{array}{ll} 48 = 1 \times 35 + 13 & 13 = 48 - 1 \times 35 \\ 35 = 2 \times 13 + 9 & 9 = 35 - 2 \times 13 \\ 13 = 1 \times 9 + 4 & 4 = 13 - 1 \times 9 \end{array}$$

Before we prove Bézout's Lemma, let's do an example. Let's take 48x - 35y = 1 that we saw earlier. Applying Euclid's algorithm, we get

| $48 = 1 \times 35 + 13$ | $13 = 48 - 1 \times 35$ |
|-------------------------|-------------------------|
| $35 = 2 \times 13 + 9$ | $9 = 35 - 2 \times 13$ |
| $13 = 1 \times 9 + 4$ | $4 = 13 - 1 \times 9$ |
| $9 = 2 \times 4 + 1$ | $1 = 9 - 2 \times 4$ |

Before we prove Bézout's Lemma, let's do an example. Let's take 48x - 35y = 1 that we saw earlier. Applying Euclid's algorithm, we get

| $48 = 1 \times 35 + 13$ | $13 = 48 - 1 \times 35$ |
|-------------------------|-------------------------|
| $35 = 2 \times 13 + 9$ | $9 = 35 - 2 \times 13$ |
| $13 = 1 \times 9 + 4$ | $4 = 13 - 1 \times 9$ |
| $9 = 2 \times 4 + 1$ | $1 = 9 - 2 \times 4$ |

$$1 = 9 + (-2) \times 4$$

Before we prove Bézout's Lemma, let's do an example. Let's take 48x - 35y = 1 that we saw earlier. Applying Euclid's algorithm, we get

| $48 = 1 \times 35 + 13$ | $13 = 48 - 1 \times 35$ |
|-------------------------|-------------------------|
| $35 = 2 \times 13 + 9$ | $9 = 35 - 2 \times 13$ |
| $13 = 1 \times 9 + 4$ | $4 = 13 - 1 \times 9$ |
| $9 = 2 \times 4 + 1$ | $1 = 9 - 2 \times 4$ |

$$1 = 9 + (-2) \times 4$$

= 9 - 2 × (13 - 1 × 9)

Before we prove Bézout's Lemma, let's do an example. Let's take 48x - 35y = 1 that we saw earlier. Applying Euclid's algorithm, we get

$$\begin{array}{ll} 48 = 1 \times 35 + 13 & 13 = 48 - 1 \times 35 \\ 35 = 2 \times 13 + 9 & 9 = 35 - 2 \times 13 \\ 13 = 1 \times 9 + 4 & 4 = 13 - 1 \times 9 \\ 9 = 2 \times 4 + 1 & 1 = 9 - 2 \times 4 \end{array}$$

$$1 = 9 + (-2) \times 4$$

= 9 - 2 × (13 - 1 × 9)
= (-2) × 13 + 3 × 9

Before we prove Bézout's Lemma, let's do an example. Let's take 48x - 35y = 1 that we saw earlier. Applying Euclid's algorithm, we get

| $48 = 1 \times 35 + 13$ | $13 = 48 - 1 \times 35$ |
|-------------------------|-------------------------|
| $35 = 2 \times 13 + 9$ | $9 = 35 - 2 \times 13$ |
| $13 = 1 \times 9 + 4$ | $4 = 13 - 1 \times 9$ |
| $9 = 2 \times 4 + 1$ | $1 = 9 - 2 \times 4$ |

$$1 = 9 + (-2) \times 4$$

= 9 - 2 × (13 - 1 × 9)
= (-2) × 13 + 3 × 9
= (-2) × 13 + 3 × (35 - 2 × 13)

Before we prove Bézout's Lemma, let's do an example. Let's take 48x - 35y = 1 that we saw earlier. Applying Euclid's algorithm, we get

| $48 = 1 \times 35 + 13$ | $13 = 48 - 1 \times 35$ |
|-------------------------|-------------------------|
| $35 = 2 \times 13 + 9$ | $9 = 35 - 2 \times 13$ |
| $13 = 1 \times 9 + 4$ | $4 = 13 - 1 \times 9$ |
| $9 = 2 \times 4 + 1$ | $1 = 9 - 2 \times 4$ |

$$1 = 9 + (-2) \times 4$$

= 9 - 2 × (13 - 1 × 9)
= (-2) × 13 + 3 × 9
= (-2) × 13 + 3 × (35 - 2 × 13)
= 3 × 35 + (-8) × 13

Before we prove Bézout's Lemma, let's do an example. Let's take 48x - 35y = 1 that we saw earlier. Applying Euclid's algorithm, we get

| $48 = 1 \times 35 + 13$ | $13 = 48 - 1 \times 35$ |
|-------------------------|-------------------------|
| $35 = 2 \times 13 + 9$ | $9 = 35 - 2 \times 13$ |
| $13 = 1 \times 9 + 4$ | $4 = 13 - 1 \times 9$ |
| $9 = 2 \times 4 + 1$ | $1 = 9 - 2 \times 4$ |

$$1 = 9 + (-2) \times 4$$

= 9 - 2 × (13 - 1 × 9)
= (-2) × 13 + 3 × 9
= (-2) × 13 + 3 × (35 - 2 × 13)
= 3 × 35 + (-8) × 13
= 3 × 35 - 8 × (48 - 1 × 35)

Before we prove Bézout's Lemma, let's do an example. Let's take 48x - 35y = 1 that we saw earlier. Applying Euclid's algorithm, we get

| $48 = 1 \times 35 + 13$ | $13 = 48 - 1 \times 35$ |
|-------------------------|-------------------------|
| $35 = 2 \times 13 + 9$ | $9 = 35 - 2 \times 13$ |
| $13 = 1 \times 9 + 4$ | $4 = 13 - 1 \times 9$ |
| $9 = 2 \times 4 + 1$ | $1 = 9 - 2 \times 4$ |

Climbing up the tower on the right-hand side,

$$1 = 9 + (-2) \times 4$$

= 9 - 2 × (13 - 1 × 9)
= (-2) × 13 + 3 × 9
= (-2) × 13 + 3 × (35 - 2 × 13)
= 3 × 35 + (-8) × 13
= 3 × 35 - 8 × (48 - 1 × 35)
= -8 × 48 + 11 × 35

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

P. E. Farrell (Oxford)

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

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

Proof.

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

$$a = q_0 b + r_0$$

 $b = q_1 r_0 + r_1$
 $r_0 = q_2 r_1 + r_2$
:

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

Let's rewrite this as

$$d = r_{i-3} - q_{i-1}r_{i-2}$$

Let's rewrite this as

$$d = r_{i-3} - q_{i-1}r_{i-2}.$$

We know that $r_{i-4} = q_{i-2}r_{i-3} + r_{i-2}$, so using this to eliminate r_{i-2} we have

$$d = -q_{i-1}r_{i-4} + (1 - q_{i-1}q_{i-2})r_{i-3}.$$

Let's rewrite this as

$$d = r_{i-3} - q_{i-1}r_{i-2}.$$

We know that $r_{i-4} = q_{i-2}r_{i-3} + r_{i-2}$, so using this to eliminate r_{i-2} we have

$$d = -q_{i-1}r_{i-4} + (1 - q_{i-1}q_{i-2})r_{i-3}.$$

Proceeding by induction, we can write d as a combination of r_{i-5} and r_{i-4} , then r_{i-6} and r_{i-5} , and so on until we write

$$d = xa + yb.$$

Let's rewrite this as

$$d = r_{i-3} - q_{i-1}r_{i-2}.$$

We know that $r_{i-4} = q_{i-2}r_{i-3} + r_{i-2}$, so using this to eliminate r_{i-2} we have

$$d = -q_{i-1}r_{i-4} + (1 - q_{i-1}q_{i-2})r_{i-3}.$$

Proceeding by induction, we can write d as a combination of r_{i-5} and r_{i-4} , then r_{i-6} and r_{i-5} , and so on until we write

$$d = xa + yb.$$

This uses an *algorithm* to prove an existence result.

P. E. Farrell (Oxford)

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

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

also. Similarly, if $a(x_p + \tilde{x}) + b(y_p + \tilde{y}) = 1$, then $a\tilde{x} + b\tilde{y} = 0$.

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

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

also. Similarly, if $a(x_p + \tilde{x}) + b(y_p + \tilde{y}) = 1$, then $a\tilde{x} + b\tilde{y} = 0$.

What are the solutions to the *homogeneous equation* $a\tilde{x} + b\tilde{y} = 0$? Exactly $(\tilde{x}, \tilde{y}) = n(-b, a)$ for $n \in \mathbb{Z}$!

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

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

also. Similarly, if $a(x_p + \tilde{x}) + b(y_p + \tilde{y}) = 1$, then $a\tilde{x} + b\tilde{y} = 0$.

What are the solutions to the homogeneous equation $a\tilde{x} + b\tilde{y} = 0$? Exactly $(\tilde{x}, \tilde{y}) = n(-b, a)$ for $n \in \mathbb{Z}$!

The general solution to ax + by = c is thus

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

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.

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.

Step 2 Divide both sides of the equation by d to get $\hat{a}x + \hat{b}y = \hat{c}$.

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.

Step 2 Divide both sides of the equation by d to get $\hat{a}x + \hat{b}y = \hat{c}$.

Step 3 Compute a *particular solution* (x_p, y_p) of $\hat{a}x + \hat{b}y = 1$.

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.
- **Step 2** Divide both sides of the equation by d to get $\hat{a}x + \hat{b}y = \hat{c}$.
- **Step 3** Compute a particular solution (x_p, y_p) of $\hat{a}x + \hat{b}y = 1$.

Step 4 Set the general solution to be

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

```
Step 1 d = \gcd(192, 140) = 4.
```

Step 1 $d = \gcd(192, 140) = 4.$

Step 2 Dividing both sides by d, we get 48x - 35y = 3.

Step 1 d = gcd(192, 140) = 4. Step 2 Dividing both sides by d, we get 48x - 35y = 3. Step 3 Solving $48x_p - 35y_p = 1$, we get $(x_p, y_p) = (-8, -11)$.

Step 1
$$d = \gcd(192, 140) = 4$$
.
Step 2 Dividing both sides by d , we get $48x - 35y = 3$.
Step 3 Solving $48x_p - 35y_p = 1$, we get $(x_p, y_p) = (-8, -11)$.
Step 4 The general solution is thus

$$\{3(-8, -11) + n(35, 48) : n \in \mathbb{Z}\}\$$

= \{(-24, -33) + n(35, 48) : n \in \mathbb{Z}\}.

Section 3

The extended Euclidean algorithm



$$ax + by = \gcd(a, b).$$

$$ax + by = \gcd(a, b).$$

The idea of climbing up the tower of equations backwards is intuitively useful, but it's not so amenable to computer implementation.

$$ax + by = \gcd(a, b).$$

The idea of climbing up the tower of equations backwards is intuitively useful, but it's not so amenable to computer implementation.

There's a very clever modification of Euclid's algorithm that computes a particular solution to the LDE in one pass: the *extended Euclidean algorithm*.

$$ax + by = \gcd(a, b).$$

The idea of climbing up the tower of equations backwards is intuitively useful, but it's not so amenable to computer implementation.

There's a very clever modification of Euclid's algorithm that computes a particular solution to the LDE in one pass: the *extended Euclidean algorithm*.

This appears to have first been explained by Āryabhaṭa (476–550).

Recall that Euclid's algorithm constructs a sequence

 $r_{-2}, r_{-1}, r_0, r_1, \ldots, r_{i-1},$

where $r_{i-1} = \text{gcd}(a, b)$ and again we denote $r_{-2} = a$, $r_{-1} = b$.

Recall that Euclid's algorithm constructs a sequence

 $r_{-2}, r_{-1}, r_0, r_1, \ldots, r_{i-1},$

where $r_{i-1} = \text{gcd}(a, b)$ and again we denote $r_{-2} = a$, $r_{-1} = b$.

We introduce two new sequences

$$x_{-2}, x_{-1}, x_0, x_1, \dots, x_{i-1}, \ y_{-2}, y_{-1}, y_0, y_1, \dots, y_{i-1},$$

Recall that Euclid's algorithm constructs a sequence

 $r_{-2}, r_{-1}, r_0, r_1, \ldots, r_{i-1},$

where $r_{i-1} = \text{gcd}(a, b)$ and again we denote $r_{-2} = a$, $r_{-1} = b$.

We introduce two new sequences

$$x_{-2}, x_{-1}, x_0, x_1, \dots, x_{i-1},$$

 $y_{-2}, y_{-1}, y_0, y_1, \dots, y_{i-1},$

and we will enforce the property that

$$ax_j + by_j = r_j, \quad j = -2, \dots, i-1.$$
Recall that Euclid's algorithm constructs a sequence

 $r_{-2}, r_{-1}, r_0, r_1, \ldots, r_{i-1},$

where $r_{i-1} = \text{gcd}(a, b)$ and again we denote $r_{-2} = a$, $r_{-1} = b$.

We introduce two new sequences

$$x_{-2}, x_{-1}, x_0, x_1, \dots, x_{i-1},$$

 $y_{-2}, y_{-1}, y_0, y_1, \dots, y_{i-1},$

and we will enforce the property that

$$ax_j + by_j = r_j, \quad j = -2, \dots, i-1.$$

If we can enforce this, then we will have

$$ax_{i-1} + by_{i-1} = r_{i-1} = \gcd(a, b).$$

How do we enforce

$$ax_j + by_j = r_j, \quad j = -2, \dots, i - 1?$$

How do we enforce

$$ax_j + by_j = r_j, \quad j = -2, \dots, i - 1?$$

Well, to begin, we should set

$$(x_{-2}, y_{-2}) = (1, 0), \quad (x_{-1}, y_{-1}) = (0, 1)$$

so that our property is enforced at the start.

How do we enforce

$$ax_j + by_j = r_j, \quad j = -2, \dots, i - 1?$$

Well, to begin, we should set

$$(x_{-2}, y_{-2}) = (1, 0), \quad (x_{-1}, y_{-1}) = (0, 1)$$

so that our property is enforced at the start.

Consider some step of Euclid's method,

$$r_j = q_{j+2}r_{j+1} + r_{j+2}.$$

If we know the expansions of r_j and r_{j+1} in terms of our 'basis' a and b, then we can work out the expansion of r_{j+2} too:

$$x_{j+2} = x_j - q_{j+2}x_{j+1},$$

$$y_{j+2} = y_j - q_{j+2}y_{j+1}.$$

Section 4

Euclid for polynomials

A *polynomial* p in $\mathbb{R}[x]$ of degree $d \in \mathbb{N}$ is an expression of the form

$$p(x) = a_0 + a_1 x + a_2 x^2 + \dots + a_d x^d$$
,

where all the a_i lie in the set of real numbers \mathbb{R} .

A *polynomial* p in $\mathbb{R}[x]$ of degree $d \in \mathbb{N}$ is an expression of the form

$$p(x) = a_0 + a_1 x + a_2 x^2 + \dots + a_d x^d$$
,

where all the a_i lie in the set of real numbers \mathbb{R} .

A root of p is a number $x \in \mathbb{C}$ satisfying p(x) = 0.

A polynomial p in $\mathbb{R}[x]$ of degree $d \in \mathbb{N}$ is an expression of the form

$$p(x) = a_0 + a_1 x + a_2 x^2 + \dots + a_d x^d$$
,

where all the a_i lie in the set of real numbers \mathbb{R} .

A root of p is a number $x \in \mathbb{C}$ satisfying p(x) = 0.

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

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

with quotient c(x) and remainder r(x), with deg(r) < deg(q).

This is an algebraic structure ${\cal R}$ that can be equipped with a $\it Euclidean$ function

$$f: R \setminus \{0\} \to \mathbb{N}$$

which is something that strictly decreases on division: given $a, b \in R$, there exist $q, r \in R$, such that

$$a = qb + r,$$

and either r = 0 or f(r) < f(b).

This is an algebraic structure ${\cal R}$ that can be equipped with a $\it Euclidean$ function

$$f: R \setminus \{0\} \to \mathbb{N}$$

which is something that strictly decreases on division: given $a, b \in R$, there exist $q, r \in R$, such that

$$a = qb + r,$$

and either r = 0 or f(r) < f(b).

For the polynomials, the Euclidean function is

$$f(r) = \deg(r).$$

This is an algebraic structure ${\cal R}$ that can be equipped with a $\it Euclidean$ function

$$f: R \setminus \{0\} \to \mathbb{N}$$

which is something that strictly decreases on division: given $a, b \in R$, there exist $q, r \in R$, such that

$$a = qb + r,$$

and either r = 0 or f(r) < f(b).

For the polynomials, the Euclidean function is

$$f(r) = \deg(r).$$

We can generalise Euclid's method, greatest common divisors, Bézout's Lemma, and many other results to such domains.

P. E. Farrell (Oxford)

For now, we focus on computing *common* roots of two polynomials p and q, of possibly different degrees. The common roots are $x \in \mathbb{C}$ such that p(x) = q(x) = 0.

For now, we focus on computing *common* roots of two polynomials p and q, of possibly different degrees. The common roots are $x \in \mathbb{C}$ such that p(x) = q(x) = 0.

We do this by finding the roots of the greatest common divisor of p and q: the polynomial of largest degree that divides both p and q.

For now, we focus on computing *common* roots of two polynomials p and q, of possibly different degrees. The common roots are $x \in \mathbb{C}$ such that p(x) = q(x) = 0.

We do this by finding the roots of the greatest common divisor of p and q: the polynomial of largest degree that divides both p and q.

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

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

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

We have

$$x^{4} + x^{3} - 6x^{2} + 5x - 1 = (x)(x^{3} + x^{2} + 3x - 5) + (-9x^{2} + 10x - 1)$$

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

We have

$$x^{4} + x^{3} - 6x^{2} + 5x - 1 = (x)(x^{3} + x^{2} + 3x - 5) + (-9x^{2} + 10x - 1)$$
$$x^{3} + x^{2} + 3x - 5 = \left(-\frac{1}{9}x - \frac{19}{81}\right)\left(-9x^{2} + 10x - 1\right) + \frac{424}{81}(x - 1)$$

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

We have

$$x^{4} + x^{3} - 6x^{2} + 5x - 1 = (x)(x^{3} + x^{2} + 3x - 5) + (-9x^{2} + 10x - 1)$$

$$x^{3} + x^{2} + 3x - 5 = \left(-\frac{1}{9}x - \frac{19}{81}\right)\left(-9x^{2} + 10x - 1\right) + \frac{424}{81}(x - 1)$$

$$-9x^{2} + 10x - 1 = -\frac{81}{424}(9x - 1)\frac{424}{81}(x - 1) + 0.$$

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

We have

$$x^{4} + x^{3} - 6x^{2} + 5x - 1 = (x)(x^{3} + x^{2} + 3x - 5) + (-9x^{2} + 10x - 1)$$

$$x^{3} + x^{2} + 3x - 5 = \left(-\frac{1}{9}x - \frac{19}{81}\right)\left(-9x^{2} + 10x - 1\right) + \frac{424}{81}(x - 1)$$

$$-9x^{2} + 10x - 1 = -\frac{81}{424}(9x - 1)\frac{424}{81}(x - 1) + 0.$$

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

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

We mention some interesting applications of Euclid's method for polynomials:

We mention some interesting applications of Euclid's method for polynomials:

1. A clever way to identify the multiple roots of a polynomial p is to compute the gcd of p and its derivative p'.

We mention some interesting applications of Euclid's method for polynomials:

1. A clever way to identify the multiple roots of a polynomial p is to compute the gcd of p and its derivative p'.

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

Many interesting polynomials are defined via recurrence relations. Euclid's method can be used to deduce facts about these without calculating them.

Many interesting polynomials are defined via recurrence relations. Euclid's method can be used to deduce facts about these without calculating them.

Consider a family of polynomials $p_k(x)$ for $k \in \mathbb{N}$ given by

$$p_0(x) = 1, \quad p_1(x) = x,$$

and

$$p_k(x) = \alpha_k(x) \times p_{k-1}(x) + \beta_k \times p_{k-2}(x),$$
 with deg $\alpha_k = 1$ and $\beta_k \in \mathbb{R} \setminus \{0\}.$

Many interesting polynomials are defined via recurrence relations. Euclid's method can be used to deduce facts about these without calculating them.

Consider a family of polynomials $p_k(x)$ for $k \in \mathbb{N}$ given by

$$p_0(x) = 1, \quad p_1(x) = x,$$

and

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

with deg $\alpha_k = 1$ and $\beta_k \in \mathbb{R} \setminus \{0\}$.

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

Chebyshev polynomials

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

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

 $T_k(x) = 2xT_{k-1}(x) - T_{k-2}(x).$



Pafnuty Chebyshev, 1821-1894

Chebyshev polynomials

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

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

 $T_k(x) = 2xT_{k-1}(x) - T_{k-2}(x).$



Pafnuty Chebyshev, 1821-1894

Laguerre polynomials

These 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



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

Proof.

First note that $\deg p_k = k$, by induction. By the recursive formula, we have

$$p_{k+1}(x) = \alpha_{k+1}(x) \times p_k(x) + \beta_{k+1} \times p_{k-1}(x).$$

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 .

Proof.

First note that $\deg p_k = k$, by induction. By the recursive formula, we have

$$p_{k+1}(x) = \alpha_{k+1}(x) \times p_k(x) + \beta_{k+1} \times p_{k-1}(x).$$

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 . Up to a nonzero scalar (which doesn't change the roots), p_{k-1} is the remainder when p_{k+1} is divided by p_k .

Proof.

First note that $\deg p_k = k$, by induction. By the recursive formula, we have

$$p_{k+1}(x) = \alpha_{k+1}(x) \times p_k(x) + \beta_{k+1} \times p_{k-1}(x).$$

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 . Up to a nonzero scalar (which doesn't change the roots), p_{k-1} is the remainder when p_{k+1} is divided by p_k .

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 $f : \mathbb{R} \to \mathbb{R}$.
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 $f : \mathbb{R} \to \mathbb{R}$.

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

Rootfinding problem

Given $f:\mathbb{R}\to\mathbb{R},$ find $x^\star\in\mathbb{R}$ such that

$$f(x^{\star}) = 0.$$

Rootfinding problem

Given $f: \mathbb{R} \to \mathbb{R}$, find $x^{\star} \in \mathbb{R}$ such that

$$f(x^{\star}) = 0.$$

This problem shows up everywhere. For example, to solve an equation

$$f_1(x) = f_2(x),$$

find a root of $f(x) \coloneqq f_1(x) - f_2(x)$.

Rootfinding problem

Given $f : \mathbb{R} \to \mathbb{R}$, find $x^* \in \mathbb{R}$ such that

$$f(x^{\star}) = 0.$$

This problem shows up everywhere. For example, to solve an equation

$$f_1(x) = f_2(x),$$

find a root of $f(x) \coloneqq f_1(x) - f_2(x)$.

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.

P. E. Farrell (Oxford)

Think back to some of the questions in Lecture 0:

- Does the algorithm terminate?
- Does the algorithm give the correct answer?
- How fast does the algorithm converge to the answer?
- How many operations does it take?

Think back to some of the questions in Lecture 0:

- Does the algorithm terminate?
- Does the algorithm give the correct answer?
- How fast does the algorithm converge to the answer?
- How many operations does it take?

Euclid's method always terminated, always gave the exact answer, and did so in a very small number of operations.

Think back to some of the questions in Lecture 0:

- Does the algorithm terminate?
- Does the algorithm give the correct answer?
- How fast does the algorithm converge to the answer?
- How many operations does it take?

Euclid's method always terminated, always gave the exact answer, and did so in a very small number of operations.

By contrast, rootfinding algorithms can only give *sequences* that converge to the root.

Think back to some of the questions in Lecture 0:

- Does the algorithm terminate?
- Does the algorithm give the correct answer?
- How fast does the algorithm converge to the answer?
- How many operations does it take?

Euclid's method always terminated, always gave the exact answer, and did so in a very small number of operations.

By contrast, rootfinding algorithms can only give *sequences* that converge to the root.

Different algorithms will trade off termination, convergence speed, and operation count.

Section 2

The first rootfinding algorithm we will meet is called the *bisection* method. It is based on the following theorem, a corollary of the Intermediate Value Theorem.

The first rootfinding algorithm we will meet is called the *bisection* method. It is based on the following theorem, a corollary of the Intermediate Value Theorem.

Bolzano's theorem (1817)

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

The statement f(a)f(b) < 0 is just a fancy way of saying f(a) and f(b) have opposite signs.



Bernhard Bolzano, 1781-1848

The first rootfinding algorithm we will meet is called the *bisection* method. It is based on the following theorem, a corollary of the Intermediate Value Theorem.

Bolzano's theorem (1817)

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

The statement f(a)f(b) < 0 is just a fancy way of saying f(a) and f(b) have opposite signs.



Bernhard Bolzano, 1781-1848

We evaluate f at c = (a + b)/2. We then have three possibilities: 1. f(c) = 0, so we are done!

The first rootfinding algorithm we will meet is called the *bisection* method. It is based on the following theorem, a corollary of the Intermediate Value Theorem.

Bolzano's theorem (1817)

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

The statement f(a)f(b) < 0 is just a fancy way of saying f(a) and f(b) have opposite signs.



Bernhard Bolzano, 1781-1848

We evaluate f at c = (a + b)/2. We then have three possibilities:

- 1. f(c) = 0, so we are done!
- 2. f(c) has the same sign as f(a), so there exists a root in (c, b).

The first rootfinding algorithm we will meet is called the *bisection* method. It is based on the following theorem, a corollary of the Intermediate Value Theorem.

Bolzano's theorem (1817)

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

The statement f(a)f(b) < 0 is just a fancy way of saying f(a) and f(b) have opposite signs.



Bernhard Bolzano, 1781-1848

We evaluate f at c = (a + b)/2. We then have three possibilities:

- 1. f(c) = 0, so we are done!
- 2. f(c) has the same sign as f(a), so there exists a root in (c, b).
- 3. f(c) has the same sign as f(b), so there exists a root in (a, c).





























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

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

function bisect(f, a, b, tol) while |b - a|/2 > tol do $c \leftarrow (a + b)/2$

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

function bisect(f, a, b, tol) while |b - a|/2 >tol do $c \leftarrow (a + b)/2$ if f(c) = 0 then return c

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

function bisect(f, a, b, tol) while |b - a|/2 >tol do $c \leftarrow (a + b)/2$ if f(c) = 0 then return celse if f(c)f(b) < 0 then a = c

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

function bisect(f, a, b, tol) while |b - a|/2 >tol do $c \leftarrow (a + b)/2$ if f(c) = 0 then return celse if f(c)f(b) < 0 then a = celse if f(a)f(c) < 0 then b = c

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

```
function bisect(f, a, b, tol)

while |b - a|/2 > tol do

c \leftarrow (a + b)/2

if f(c) = 0 then return c

else if f(c)f(b) < 0 then a = c

else if f(a)f(c) < 0 then b = c

end if

end while
```

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

```
function bisect(f, a, b, tol)

while |b - a|/2 > \text{tol do}

c \leftarrow (a + b)/2

if f(c) = 0 then return c

else if f(c)f(b) < 0 then a = c

else if f(a)f(c) < 0 then b = c

end if

end while

return (a + b)/2

end function
```
Let's state this as an algorithm.

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

```
function bisect(f, a, b, tol)

while |b - a|/2 > \text{tol do}

c \leftarrow (a + b)/2

if f(c) = 0 then return c

else if f(c)f(b) < 0 then a = c

else if f(a)f(c) < 0 then b = c

end if

end while

return (a + b)/2

end function
```

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

P. E. Farrell (Oxford)

Bisection

There's not much published information on the history of bisection. The earliest reference Prof. Hollings could find to it was in Cauchy's *Cours d'analyse* (1821).



Augustin-Louis Cauchy FRS 1789–1857

Bisection

There's not much published information on the history of bisection. The earliest reference Prof. Hollings could find to it was in Cauchy's *Cours d'analyse* (1821).



Augustin-Louis Cauchy FRS 1789–1857

Lemma

The algorithm always terminates.

Bisection

There's not much published information on the history of bisection. The earliest reference Prof. Hollings could find to it was in Cauchy's *Cours d'analyse* (1821).



Augustin-Louis Cauchy FRS 1789–1857

Lemma

The algorithm always terminates.

Proof.

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

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

Let's start with $[a,b]=[-10,10].~f(-10)\approx-9.16,~f(10)\approx10.83,$ so we're good to go.

| c | f(c) | [a,b] |
|---|------|---------|
| 0 | -1 | [0, 10] |

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

| с | f(c) | [a,b] |
|---|------|---------|
| 0 | -1 | [0, 10] |
| 5 | 4.71 | [0, 5] |

Let's start with $[a,b]=[-10,10].~f(-10)\approx-9.16,~f(10)\approx10.83,$ so we're good to go.

| с | f(c) | [a,b] |
|-----|------|----------|
| 0 | -1 | [0, 10] |
| 5 | 4.71 | [0, 5] |
| 2.5 | 3.30 | [0, 2.5] |

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

| с | f(c) | [a,b] |
|------|------|-----------|
| 0 | -1 | [0, 10] |
| 5 | 4.71 | [0, 5] |
| 2.5 | 3.30 | [0, 2.5] |
| 1.25 | 0.93 | [0, 1.25] |

Let's start with $[a,b]=[-10,10].~f(-10)\approx-9.16,~f(10)\approx10.83,$ so we're good to go.

| с | f(c) | [a,b] |
|----------------------|------------------------|--|
| 0 | -1 | [0, 10] |
| 5 | 4.71 | [0, 5] |
| 2.5 | 3.30 | [0, 2.5] |
| 1.25 | 0.93 | [0, 1.25] |
| 0.625 | -0.185 | [0.625, 1.25] |
| 2.5 1.25 0.625 | 3.30 0.93 -0.185 | [0, 2.5] [0, 1.25] [0.625, 1.25] |

Let's start with $[a,b]=[-10,10].~f(-10)\approx-9.16,~f(10)\approx10.83,$ so we're good to go.

| с | f(c) | [a,b] |
|-------|--------|---------------|
| 0 | -1 | [0, 10] |
| 5 | 4.71 | [0, 5] |
| 2.5 | 3.30 | [0, 2.5] |
| 1.25 | 0.93 | [0, 1.25] |
| 0.625 | -0.185 | [0.625, 1.25] |

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

✓ When it applies, it is guaranteed to converge.

- \checkmark When it applies, it is guaranteed to converge.
- The method is very simple and very robust.

- \checkmark When it applies, it is guaranteed to converge.
- ✓ The method is very simple and very robust.
- ✓ Its smoothness requirements are low (only continuity).

- \checkmark When it applies, it is guaranteed to converge.
- ✓ The method is very simple and very robust.
- ✓ Its smoothness requirements are low (only continuity).
- X The interval of interest only reduces by a factor of two each time.

- \checkmark When it applies, it is guaranteed to converge.
- ✓ The method is very simple and very robust.
- ✓ Its smoothness requirements are low (only continuity).
- X The interval of interest only reduces by a factor of two each time.
- \checkmark It can be hard to find the initial points [a, b] that bracket a root.

- \checkmark When it applies, it is guaranteed to converge.
- The method is very simple and very robust.
- ✓ Its smoothness requirements are low (only continuity).
- X The interval of interest only reduces by a factor of two each time.
- \checkmark It can be hard to find the initial points [a, b] that bracket a root.
- × It is hard (but not impossible) to generalise to higher dimensions.

- \checkmark When it applies, it is guaranteed to converge.
- The method is very simple and very robust.
- ✓ Its smoothness requirements are low (only continuity).
- X The interval of interest only reduces by a factor of two each time.
- \checkmark It can be hard to find the initial points [a, b] that bracket a root.
- × It is hard (but not impossible) to generalise to higher dimensions.
- X It can never find roots of even multiplicity.

- \checkmark When it applies, it is guaranteed to converge.
- The method is very simple and very robust.
- ✓ Its smoothness requirements are low (only continuity).
- X The interval of interest only reduces by a factor of two each time.
- \checkmark It can be hard to find the initial points [a, b] that bracket a root.
- × It is hard (but not impossible) to generalise to higher dimensions.
- X It can never find roots of even multiplicity.

Definition (Multiplicity of a root)

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

- \checkmark When it applies, it is guaranteed to converge.
- The method is very simple and very robust.
- ✓ Its smoothness requirements are low (only continuity).
- X The interval of interest only reduces by a factor of two each time.
- \checkmark It can be hard to find the initial points [a, b] that bracket a root.
- × It is hard (but not impossible) to generalise to higher dimensions.
- X It can never find roots of even multiplicity.

Definition (Multiplicity of a root)

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

Later we will study other methods with different sets of advantages and disadvantages.

P. E. Farrell (Oxford)

Section 3

Rate of convergence of a sequence

Definition (Linear convergence of a sequence)

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

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

Definition (Linear convergence of a sequence)

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

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

In other words, asymptotically, moving one step along the sequence multiplies the error by a fixed $\mu < 1$. The μ is called the *rate of convergence*.

Definition (Linear convergence of a sequence)

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

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

In other words, asymptotically, moving one step along the sequence multiplies the error by a fixed $\mu < 1$. The μ is called the *rate of convergence*.

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

Definition (Superlinear convergence of a sequence)

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

$$\lim_{k \to \infty} \frac{|x_{i+1} - x^*|}{|x_i - x^*|} = 0.$$

Definition (Superlinear convergence of a sequence)

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

$$\lim_{k \to \infty} \frac{|x_{i+1} - x^*|}{|x_i - x^*|} = 0.$$

In other words, the sequence converges faster than any linear rate of convergence.

Definition (Superlinear convergence of a sequence)

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

$$\lim_{k \to \infty} \frac{|x_{i+1} - x^*|}{|x_i - x^*|} = 0.$$

In other words, the sequence converges faster than any linear rate of convergence.

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.

P. E. Farrell (Oxford)

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

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

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

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

We can further classify superlinear convergence:

Definition (Order of convergence of a sequence)

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

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

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

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

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

It will be very useful to consider finding *fixed points*: given $g:[a,b] \to \mathbb{R}$, find $x^* \in [a,b]$ such that $g(x^*) = x^*$.

It will be very useful to consider finding *fixed points*: given $g : [a, b] \to \mathbb{R}$, find $x^* \in [a, b]$ such that $g(x^*) = x^*$.

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.

It will be very useful to consider finding *fixed points*: given $g : [a, b] \to \mathbb{R}$, find $x^* \in [a, b]$ such that $g(x^*) = x^*$.

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.

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.
So far we have considered rootfinding: find $x^{\star} \in \mathbb{R}$ such that $f(x^{\star}) = 0$.

It will be very useful to consider finding *fixed points*: given $g : [a, b] \to \mathbb{R}$, find $x^* \in [a, b]$ such that $g(x^*) = x^*$.

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.

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



Luitzen Brouwer, 1881-1966

When can we show fixed points exist?

Theorem (Brouwer's fixed point theorem)

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



```
Luitzen Brouwer, 1881-1966
```

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.

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 $f(b) \leq 0$.

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 $f(b) \leq 0$.

If either inequality is an equality, we have a fixed point. So assume that f(a)>0 and $f(b)<0. \label{eq:factor}$

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 $f(b) \leq 0$.

If either inequality is an equality, we have a fixed point. So assume that f(a)>0 and $f(b)<0. \label{eq:factor}$

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.

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

Theorem

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

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

Theorem

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

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

There must be at least one fixed point of g, since it is continuous.

There must be at least one fixed point of g, since it is continuous.

Suppose p and q are two fixed points of g in (a, b), then we have

$$g(p) = p, \quad g(q) = q.$$

There must be at least one fixed point of g, since it is continuous.

Suppose p and q are two fixed points of g in (a, b), then we have

$$g(p) = p, \quad g(q) = q.$$

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

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

There must be at least one fixed point of g, since it is continuous.

Suppose p and q are two fixed points of g in (a, b), then we have

$$g(p) = p, \quad g(q) = q.$$

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

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

But |g'(r)| < 1 by assumption, a contradiction.

There must be at least one fixed point of g, since it is continuous.

Suppose p and q are two fixed points of g in (a, b), then we have

$$g(p) = p, \quad g(q) = q.$$

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

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

But |g'(r)| < 1 by assumption, a contradiction.

How do we turn this into an algorithm?

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

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

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

```
function fixedpoint(g, x_0, tol)

x \leftarrow x_0

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

x \leftarrow g(x)

end while
```

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

```
function fixedpoint(g, x_0, tol)

x \leftarrow x_0

while |g(x) - x| > tol do

x \leftarrow g(x)

end while

return g(x)

end function
```

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

```
function fixedpoint(g, x_0, tol)

x \leftarrow x_0

while |g(x) - x| > tol do

x \leftarrow g(x)

end while

return g(x)

end function
```

Our goal is to investigate when this converges.

















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^* .

We then proposed the iteration scheme: take any $x_0 \in [a, b]$, and set

$$x_{i+1} = g(x_i)$$

until convergence.

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^* .

We then proposed the iteration scheme: take any $x_0 \in [a, b]$, and set

$$x_{i+1} = g(x_i)$$

until convergence.

Thus, we are checking $x_0, g(x_0), g(g(x_0)), \ldots$ to find the unique fixed point.

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^* .

We then proposed the iteration scheme: take any $x_0 \in [a, b]$, and set

$$x_{i+1} = g(x_i)$$

until convergence.

Thus, we are checking $x_0, g(x_0), g(g(x_0)), \ldots$ to find the unique fixed point.

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

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

for all $x, y \in [a, b]$.

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

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

for all $x, y \in [a, b]$.

Example

Any differentiable $g : [a, b] \rightarrow [a, b]$ with $|g'(x)| \le \gamma < 1$ 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

$$|g(x) - g(y)| = |g'(c)(x - y)| \le \gamma |x - y|.$$

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

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

for all $x, y \in [a, b]$.

Example

Any differentiable $g : [a, b] \to [a, b]$ with $|g'(x)| \le \gamma < 1$ 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

$$|g(x) - g(y)| = |g'(c)(x - y)| \le \gamma |x - y|.$$

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^* , and the iteration scheme $x_{i+1} = g(x_i)$ converges at least linearly to x^* for any $x_0 \in [a, b]$.

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



Stefan Banach, 1892-1945

Contraction mapping theorem (1922)

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

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

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



Stefan Banach, 1892-1945

We prove the theorem in stages. First, we show g is continuous, and thus must have a fixed point.
We prove the theorem in stages. First, we show g is continuous, and thus must have a fixed point.

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

We prove the theorem in stages. First, we show g is continuous, and thus must have a fixed point.

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

We thus know that g must have a fixed point.

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

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

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

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

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

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

$$|x_i - x^*| = |g(x_{i-1}) - g(x^*)| \le \gamma |x_{i-1} - x^*| \le \gamma^2 |x_{i-2} - x^*|$$

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

$$|x_{i} - x^{\star}| = |g(x_{i-1}) - g(x^{\star})| \le \gamma |x_{i-1} - x^{\star}|$$
$$\le \gamma^{2} |x_{i-2} - x^{\star}|$$
$$\le \gamma^{i} |x_{0} - x^{\star}|.$$

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

$$|x_{i} - x^{\star}| = |g(x_{i-1}) - g(x^{\star})| \le \gamma |x_{i-1} - x^{\star}|$$
$$\le \gamma^{2} |x_{i-2} - x^{\star}|$$
$$\le \gamma^{i} |x_{0} - x^{\star}|.$$

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

$$\lim_{i \to \infty} |x_i - x^\star| = 0,$$

i.e. $x_i \to x^{\star}$.

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

$$|x_{i} - x^{\star}| = |g(x_{i-1}) - g(x^{\star})| \le \gamma |x_{i-1} - x^{\star}|$$
$$\le \gamma^{2} |x_{i-2} - x^{\star}|$$
$$\le \gamma^{i} |x_{0} - x^{\star}|.$$

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

$$\lim_{i \to \infty} |x_i - x^\star| = 0,$$

i.e. $x_i \to x^*$. Since

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

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

P. E. Farrell (Oxford)

Existence of fixed point: $g:[a,b] \rightarrow [a,b]$ continuous.

Existence of fixed point: $g:[a,b] \rightarrow [a,b]$ continuous.

Uniqueness of fixed point: g differentiable with |g'(x)| < 1 for all $x \in (a,b).$

Existence of fixed point: $g:[a,b] \rightarrow [a,b]$ continuous.

Uniqueness of fixed point: g differentiable with |g'(x)| < 1 for all $x \in (a,b).$

Contraction mapping theorem: g a contraction, i.e. with $|g'(x)| \leq \gamma < 1$ for all $x \in (a,b)$ in the differentiable case.

Existence of fixed point: $g:[a,b] \rightarrow [a,b]$ continuous.

Uniqueness of fixed point: g differentiable with |g'(x)| < 1 for all $x \in (a,b).$

Contraction mapping theorem: g a contraction, i.e. with $|g'(x)| \leq \gamma < 1$ for all $x \in (a,b)$ in the differentiable case.

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.

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.

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

Section 6

Example

Suppose we wish to find the roots of $f(x) = x^2 - x - 1 = 0$. (Its roots are the golden ratio $\phi \approx 1.61834$ and its conjugate $-\phi^{-1} \approx -0.618034$.)

Suppose we wish to find the roots of $f(x) = x^2 - x - 1 = 0$. (Its roots are the golden ratio $\phi \approx 1.61834$ and its conjugate $-\phi^{-1} \approx -0.618034$.)

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

Fixed point iteration A

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

Suppose we wish to find the roots of $f(x) = x^2 - x - 1 = 0$. (Its roots are the golden ratio $\phi \approx 1.61834$ and its conjugate $-\phi^{-1} \approx -0.618034$.)

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

Fixed point iteration A

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

Suppose we wish to find the roots of $f(x) = x^2 - x - 1 = 0$. (Its roots are the golden ratio $\phi \approx 1.61834$ and its conjugate $-\phi^{-1} \approx -0.618034$.)

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

Fixed point iteration A

$$x^2 - x - 1 = 0 \implies x^2 = x + 1 \implies x = (x + 1)/x \Rightarrow = g_A(x)$$

Suppose we wish to find the roots of $f(x) = x^2 - x - 1 = 0$. (Its roots are the golden ratio $\phi \approx 1.61834$ and its conjugate $-\phi^{-1} \approx -0.618034$.)

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

Fixed point iteration A

$$x^2 - x - 1 = 0 \implies x^2 = x + 1 \implies x = (x + 1)/x \Rightarrow g_A(x)$$

Fixed point iteration B

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

Suppose we wish to find the roots of $f(x) = x^2 - x - 1 = 0$. (Its roots are the golden ratio $\phi \approx 1.61834$ and its conjugate $-\phi^{-1} \approx -0.618034$.)

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

Fixed point iteration A

$$x^2 - x - 1 = 0 \implies x^2 = x + 1 \implies x = (x + 1)/x \Rightarrow = g_A(x)$$

Fixed point iteration B

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

Suppose we wish to find the roots of $f(x) = x^2 - x - 1 = 0$. (Its roots are the golden ratio $\phi \approx 1.61834$ and its conjugate $-\phi^{-1} \approx -0.618034$.)

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

Fixed point iteration A

$$x^2 - x - 1 = 0 \implies x^2 = x + 1 \implies x = (x + 1)/x \Rightarrow = g_A(x)$$

Fixed point iteration B

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

Fixed point iteration C

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

P. E. Farrell (Oxford)

Suppose we wish to find the roots of $f(x) = x^2 - x - 1 = 0$. (Its roots are the golden ratio $\phi \approx 1.61834$ and its conjugate $-\phi^{-1} \approx -0.618034$.)

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

Fixed point iteration A

$$x^2 - x - 1 = 0 \implies x^2 = x + 1 \implies x = (x + 1)/x \Rightarrow = g_A(x)$$

Fixed point iteration B

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

Fixed point iteration C

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

P. E. Farrell (Oxford)

Suppose we wish to find the roots of $f(x) = x^2 - x - 1 = 0$. (Its roots are the golden ratio $\phi \approx 1.61834$ and its conjugate $-\phi^{-1} \approx -0.618034$.)

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

Fixed point iteration A

$$x^2 - x - 1 = 0 \implies x^2 = x + 1 \implies x = (x + 1)/x \Rightarrow = g_A(x)$$

Fixed point iteration B

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

Fixed point iteration C

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

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 1 | the | fixed | point | iteration | with | $x_0 =$ | 1.1, | we ge | et |
|-------|-------|-----|-------|-------|-----------|------|---------|------|-------|----|
|-------|-------|-----|-------|-------|-----------|------|---------|------|-------|----|

| 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 |
| 9 | 1.618151 | -0.000000 | -0.619318 |
| 10 | 1.617989 | -1.000000 | -0.617544 |

| If we r | run the | fixed | point | iteration | with x_0 | $_{0} = 1.1$ | 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 |
| 9 | 1.618151 | -0.000000 | -0.619318 |
| 10 | 1.617989 | -1.000000 | -0.617544 |

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 $\gamma = |g'(1.1)| \approx 0.826 < 1$.

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 $\gamma = |g'(1.1)| \approx 0.826 < 1$.

We also need to check that $g([a,b]) \subset [a,b]$. g(x) = 1 + 1/x, so the function is decreasing on [a,b]. Checking, we find g(1.1) = 1.9 and g(2) = 1.5, so this is satisfied.

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 $\gamma = |g'(1.1)| \approx 0.826 < 1$.

We also need to check that $g([a,b]) \subset [a,b]$. g(x) = 1 + 1/x, so the function is decreasing on [a,b]. Checking, we find g(1.1) = 1.9 and g(2) = 1.5, so this is satisfied.

Banach's contraction mapping theorem thus applies.

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 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 $\gamma = |g'(-0.4)| \approx 0.51$.

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 $\gamma = |g'(-0.4)| \approx 0.51$.

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 $g([a,b]) \subset [a,b]$.

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 $\gamma = |g'(-0.4)| \approx 0.51$.

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 $g([a,b]) \subset [a,b]$.

Banach's contraction mapping theorem thus applies.

Section 7

Termination criteria
In the proof, we saw that $|x_i - x^*| \leq \gamma^i |x_0 - x^*|$. Since $x_0, x^* \in [a, b]$, we can bound this by $\gamma^i |b - a|$.

In the proof, we saw that $|x_i - x^*| \leq \gamma^i |x_0 - x^*|$. Since $x_0, x^* \in [a, b]$, we can bound this by $\gamma^i |b - a|$.

Thus, to achieve a tolerance tol on the error, we choose i such that $\gamma^i \leq {\rm tol}/|b-a|.$

In the proof, we saw that $|x_i - x^*| \leq \gamma^i |x_0 - x^*|$. Since $x_0, x^* \in [a, b]$, we can bound this by $\gamma^i |b - a|$.

Thus, to achieve a tolerance tol on the error, we choose i such that $\gamma^i \leq {\rm tol}/|b-a|.$

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

In the proof, we saw that $|x_i - x^*| \leq \gamma^i |x_0 - x^*|$. Since $x_0, x^* \in [a, b]$, we can bound this by $\gamma^i |b - a|$.

Thus, to achieve a tolerance tol on the error, we choose i such that $\gamma^i \leq {\rm tol}/|b-a|.$

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?

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

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

$$|x_J - x_i| = |(x_J - x_{J-1}) + (x_{J-1} - x_{J-2}) + \dots + (x_{i+1} - x_i)|$$

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

$$|x_J - x_i| = |(x_J - x_{J-1}) + (x_{J-1} - x_{J-2}) + \dots + (x_{i+1} - x_i)|$$

$$\leq |x_J - x_{J-1}| + |x_{J-1} - x_{J-2}| + \dots + |x_{i+1} - x_i|$$

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

$$\begin{aligned} |x_J - x_i| &= |(x_J - x_{J-1}) + (x_{J-1} - x_{J-2}) + \dots + (x_{i+1} - x_i)| \\ &\leq |x_J - x_{J-1}| + |x_{J-1} - x_{J-2}| + \dots + |x_{i+1} - x_i| \\ &\leq \gamma^{J-1} |x_1 - x_0| + \gamma^{J-2} |x_1 - x_0| + \dots + \gamma^i |x_1 - x_0| \end{aligned}$$

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

$$\begin{aligned} |x_J - x_i| &= |(x_J - x_{J-1}) + (x_{J-1} - x_{J-2}) + \dots + (x_{i+1} - x_i)| \\ &\leq |x_J - x_{J-1}| + |x_{J-1} - x_{J-2}| + \dots + |x_{i+1} - x_i| \\ &\leq \gamma^{J-1} |x_1 - x_0| + \gamma^{J-2} |x_1 - x_0| + \dots + \gamma^i |x_1 - x_0| \\ &= (\gamma^{J-1} + \gamma^{J-2} + \dots \gamma^i) |x_1 - x_0| \end{aligned}$$

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

$$\begin{aligned} |x_J - x_i| &= |(x_J - x_{J-1}) + (x_{J-1} - x_{J-2}) + \dots + (x_{i+1} - x_i)| \\ &\leq |x_J - x_{J-1}| + |x_{J-1} - x_{J-2}| + \dots + |x_{i+1} - x_i| \\ &\leq \gamma^{J-1} |x_1 - x_0| + \gamma^{J-2} |x_1 - x_0| + \dots + \gamma^i |x_1 - x_0| \\ &= (\gamma^{J-1} + \gamma^{J-2} + \dots + \gamma^i) |x_1 - x_0| \\ &= \gamma^i \left(\gamma^{J-i-1} + \gamma^{J-i-2} + \dots + \gamma + 1 \right) |x_1 - x_0|. \end{aligned}$$

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

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

$$\begin{aligned} |x_J - x_i| &= |(x_J - x_{J-1}) + (x_{J-1} - x_{J-2}) + \dots + (x_{i+1} - x_i)| \\ &\leq |x_J - x_{J-1}| + |x_{J-1} - x_{J-2}| + \dots + |x_{i+1} - x_i| \\ &\leq \gamma^{J-1} |x_1 - x_0| + \gamma^{J-2} |x_1 - x_0| + \dots + \gamma^i |x_1 - x_0| \\ &= (\gamma^{J-1} + \gamma^{J-2} + \dots + \gamma^i) |x_1 - x_0| \\ &= \gamma^i \left(\gamma^{J-i-1} + \gamma^{J-i-2} + \dots + \gamma + 1 \right) |x_1 - x_0|. \end{aligned}$$

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

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

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

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

$$\begin{aligned} |x_J - x_i| &= |(x_J - x_{J-1}) + (x_{J-1} - x_{J-2}) + \dots + (x_{i+1} - x_i)| \\ &\leq |x_J - x_{J-1}| + |x_{J-1} - x_{J-2}| + \dots + |x_{i+1} - x_i| \\ &\leq \gamma^{J-1} |x_1 - x_0| + \gamma^{J-2} |x_1 - x_0| + \dots + \gamma^i |x_1 - x_0| \\ &= (\gamma^{J-1} + \gamma^{J-2} + \dots \gamma^i) |x_1 - x_0| \\ &= \gamma^i \left(\gamma^{J-i-1} + \gamma^{J-i-2} + \dots + \gamma + 1\right) |x_1 - x_0|. \end{aligned}$$

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

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

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

P. E. Farrell (Oxford)

Section 8

Another example

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

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

We need:

(i) $g:[a,b] \rightarrow [a,b]$, and (ii) $|g'(x)| \leq \gamma < 1$ on [a,b] for some γ .

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

We need:

(i) $g: [a, b] \rightarrow [a, b]$, and (ii) $|g'(x)| \leq \gamma < 1$ on [a, b] for some γ .

So let's consider g'(x). Calculating, we find $g'(x) = -e^{-x}$, so $|g'(x)| = |e^{-x}|$. This is 1 at x = 0 and strictly less than 1 for x > 0.

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

We need:

(i) $g:[a,b] \rightarrow [a,b]$, and (ii) $|g'(x)| \leq \gamma < 1$ on [a,b] for some γ .

So let's consider g'(x). Calculating, we find $g'(x) = -e^{-x}$, so $|g'(x)| = |e^{-x}|$. This is 1 at x = 0 and strictly less than 1 for x > 0.

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

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

We need:

(i) $g:[a,b] \rightarrow [a,b]$, and (ii) $|g'(x)| \leq \gamma < 1$ on [a,b] for some γ .

So let's consider g'(x). Calculating, we find $g'(x) = -e^{-x}$, so $|g'(x)| = |e^{-x}|$. This is 1 at x = 0 and strictly less than 1 for x > 0.

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 $x^* \approx 0.567143$.)

Our γ is $e^{-1/10} \approx 0.905$.

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.)

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 $tol = 10^{-3}$ and i > 116 for $tol = 10^{-6}$.

Section 9

Bonus: accelerating sequence convergence

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

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

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

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

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

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

Equating the two expressions for μ and doing some algebra yields

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

so we hope that the expression on the right gives a good approximation to the sequence limit.

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

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

Equating the two expressions for μ and doing some algebra yields

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

so we hope that the expression on the right gives a good approximation to the sequence limit.

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

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

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

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

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.

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

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

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.

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:

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

Let's consider rootfinding again:

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

Since there are powerful theorems about fixed point problems, let's try to reformulate this as a fixed point problem:

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

Let's consider rootfinding again:

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

Since there are powerful theorems about fixed point problems, let's try to reformulate this as a fixed point problem:

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

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



Build a *linear model* of the function.

Newton's method

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



Set x_1 to be the root of the linear model.



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



Repeat.

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

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

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

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

and solving for x_{i+1} yields

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

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

This is a generic way of constructing a fixed point problem x = g(x) from a rootfinding problem f(x) = 0. Is:



Isaac Newton FRS, 1643-1727

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

This is a generic way of constructing a fixed point



problem x = q(x) from a rootfinding problem f(x) = 0. Isaac Newton FRS, 1643–1727

The special case of applying Newton's method for calculating square roots was known to the ancient Greeks in Alexandria (Heron's method, 60).

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

This is a generic way of constructing a fixed point problem x = q(x) from a rootfinding problem f(x) = 0.



The special case of applying Newton's method for calculating square roots was known to the ancient Greeks in Alexandria (Heron's method, 60).

Taking
$$f(x) = x^2 - c$$
, we get

$$x_{i+1} = x_i - \frac{x_i^2 - c}{2x_i}$$

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





Isaac Newton FRS, 1643–1727

The special case of applying Newton's method for calculating square roots was known to the ancient Greeks in Alexandria (Heron's method, 60).

Taking
$$f(x) = x^2 - c$$
, we get

$$x_{i+1} = x_i - \frac{x_i^2 - c}{2x_i} = \frac{1}{2} \left(x_i + \frac{c}{x_i} \right).$$

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

This is a generic way of constructing a fixed point problem x = q(x) from a rootfinding problem f(x) = 0.



problem x = g(x) from a rootfinding problem f(x) = 0. Isaac Newton FRS, 1643–1727

The special case of applying Newton's method for calculating square roots was known to the ancient Greeks in Alexandria (Heron's method, 60).

Taking
$$f(x) = x^2 - c$$
, we get

$$x_{i+1} = x_i - \frac{x_i^2 - c}{2x_i} = \frac{1}{2} \left(x_i + \frac{c}{x_i} \right).$$

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

P. E. Farrell (Oxford)

John Wallis (Savilian Chair of Geometry in Oxford) published the same method before Newton, in 1685. So we should probably call it the Wallis method!



John Wallis, 1616-1703

John Wallis (Savilian Chair of Geometry in Oxford) published the same method before Newton, in 1685. So we should probably call it the Wallis method!



John Wallis, 1616-1703

Joseph Raphson (1690) simplified the method, but still only applied it to polynomials.

John Wallis (Savilian Chair of Geometry in Oxford) published the same method before Newton, in 1685. So we should probably call it the Wallis method!



John Wallis, 1616-1703

Joseph Raphson (1690) simplified the method, but still only applied it to polynomials.

Thomas Simpson (1740) gave the modern description, using calculus, and applied it to general functions.

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

Comments:

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

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

- ✓ If $f(x_i) = 0$, then $x_{i+1} = x_i$. So roots of f are fixed points of g.
- \checkmark Unlike bisection, we require f to be differentiable.

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

- ✓ If $f(x_i) = 0$, then $x_{i+1} = x_i$. So roots of f are fixed points of g.
- \times Unlike bisection, we require f to be differentiable.
- × Moreover, we need $f'(x_i) \neq 0$ at every iterate.

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

- ✓ If $f(x_i) = 0$, then $x_{i+1} = x_i$. So roots of f are fixed points of g.
- \checkmark Unlike bisection, we require f to be differentiable.
- × Moreover, we need $f'(x_i) \neq 0$ at every iterate.
- ✓ If x_0 is close to x^* , Newton's method usually converges very fast.

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

- ✓ If $f(x_i) = 0$, then $x_{i+1} = x_i$. So roots of f are fixed points of g.
- \checkmark Unlike bisection, we require f to be differentiable.
- × Moreover, we need $f'(x_i) \neq 0$ at every iterate.
- ✓ If x_0 is close to x^* , Newton's method usually converges very fast.
- \checkmark If x_0 is far away, the method can diverge or get stuck in a cycle.

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

- ✓ If $f(x_i) = 0$, then $x_{i+1} = x_i$. So roots of f are fixed points of g.
- \checkmark Unlike bisection, we require f to be differentiable.
- × Moreover, we need $f'(x_i) \neq 0$ at every iterate.
- ✓ If x_0 is close to x^* , Newton's method usually converges very fast.
- × If x_0 is far away, the method can diverge or get stuck in a cycle.
- ✓ Newton's method generalises elegantly to higher dimensions.







Even when it converges, Newton's method can behave in an unstable manner.

Even when it converges, Newton's method can behave in an unstable manner.

$$f(x) = (x-4)(x-1)(x+3), \quad x_0 = 2.352836327.$$

Even when it converges, Newton's method can behave in an unstable manner.

$$f(x) = (x-4)(x-1)(x+3), \quad x_0 = 2.352836327.$$

```
In [14]: newton(lambda x: (x-4)*(x-1)*(x+3),
                lambda x: 3*x**2 - 4*x - 11, 2.352836327, 1e-6)
Iteration 0: x = 2.352836e+00 f(x) = -1.192795e+01
Iteration 1: x = -7.829394e-01 f(x) = 1.890641e+01
Iteration 2: x = 2.352836e+00 f(x) = -1.192796e+01
Iteration 3: x = -7.829406e-01 f(x) = 1.890641e+01
. . .
Iteration 9: x = -8.476712e-01 f(x) = 1.927820e+01
Iteration 10: x = 2.687229e+00 f(x) = -1.259690e+01
Iteration 11: x = -1.449560e+02 f(x) = -3.086271e+06
Iteration 12: x = -9.643403e+01 f(x) = -9.143167e+05
. . .
Iteration 19: x = -5.622219e+00 f(x) = -1.670889e+02
Iteration 20: x = -4.050607e+00 f(x) = -4.271814e+01
Iteration 21: x = -3.265703e+00 f(x) = -8.235014e+00
Iteration 22: x = -3.023904e+00 f(x) = -6.756020e-01
Iteration 23: x = -3.000221e+00 f(x) = -6.196356e-03
Iteration 24: x = -3.000000e+00 f(x) = -5.385373e-07
D_{11}+ [14] : -3,0000000192334735
```

P. E. Farrell (Oxford)

Now change from $x_0 = 2.352836327$ to $x_0 = 2.352836323$.

Now change from $x_0 = 2.352836327$ to $x_0 = 2.352836323$.

In [15]: newton(lambda x: (x-4)*(x-1)*(x+3), lambda x: 3*x**2 - 4*x - 11, 2.352836323, 1e-6) Iteration 0: x = 2.352836e+00 f(x) = -1.192795e+01Iteration 1: x = -7.829394e-01 f(x) = 1.890641e+01Iteration 2: x = 2.352836e+00 f(x) = -1.192795e+01Iteration 3: x = -7.829393e-01 f(x) = 1.890641e+01Iteration 4: x = 2.352836e+00 f(x) = -1.192795e+01Iteration 5: x = -7.829361e-01 f(x) = 1.890639e+01Iteration 6: x = 2.352822e+00 f(x) = -1.192790e+01Iteration 7: x = -7.828166e-01 f(x) = 1.890567e+01Iteration 8: x = 2.352281e+00 f(x) = -1.192584e+01Iteration 9: x = -7.783146e-01 f(x) = 1.887843e+01Iteration 10: x = 2.332103e+00 f(x) = -1.184692e+01Iteration 11: x = -6.205467e-01 f(x) = 1.781690e+01Iteration 12: x = 1.799380e+00 f(x) = -8.442739e+00Iteration 13: x = 8.042685e-01 f(x) = 2.379590e+00Iteration 14: x = 9.981010e-01 f(x) = 2.279200e-02Iteration 15: x = 9.999997e-01 f(x) = 3.591499e-06Iteration 16: x = 1.000000e+00 f(x) = 8.926193e-14Out[15]: 0.9999999999999926

P. E. Farrell (Oxford)













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^* for a moment.

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^* for a moment.

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

$$|x_{i+1} - x^\star| \le \gamma |x_i - x^\star|,$$

or in other words that we have linear convergence

$$\frac{|x_{i+1} - x^*|}{|x_i - x^*|} \le \gamma < 1.$$
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^* for a moment.

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

$$|x_{i+1} - x^\star| \le \gamma |x_i - x^\star|,$$

or in other words that we have linear convergence

$$\frac{|x_{i+1} - x^*|}{|x_i - x^*|} \le \gamma < 1.$$

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

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

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^* for a moment.

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

$$|x_{i+1} - x^\star| \le \gamma |x_i - x^\star|,$$

or in other words that we have linear convergence

$$\frac{|x_{i+1} - x^*|}{|x_i - x^*|} \le \gamma < 1.$$

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

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

Recall that we called this *quadratic* convergence.

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

P. E. Farrell (Oxford)

Recall the Taylor expansion of g around some point a:

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

Recall the Taylor expansion of g around some point a:

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

What happens if we evaluate this around a fixed point x^* of g?

Recall the Taylor expansion of g around some point a:

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

What happens if we evaluate this around a fixed point x^* of g?

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

Recall the Taylor expansion of g around some point a:

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

What happens if we evaluate this around a fixed point x^* of g?

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

But $g(x_i) = x_{i+1}$ and $g(x^{\star}) = x^{\star}$, so

$$|x_{i+1} - x^{\star}| = |(x_i - x^{\star})g'(x^{\star}) + \frac{1}{2}(x_i - x^{\star})^2 g''(\zeta_i)|$$

Recall the Taylor expansion of g around some point a:

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

What happens if we evaluate this around a fixed point x^* of g?

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

But $g(x_i) = x_{i+1}$ and $g(x^{\star}) = x^{\star}$, so

$$|x_{i+1} - x^{\star}| = |(x_i - x^{\star})g'(x^{\star}) + \frac{1}{2}(x_i - x^{\star})^2 g''(\zeta_i)|$$

$$\leq |x_i - x^{\star}||g'(x^{\star})| + \frac{1}{2}|x_i - x_{\star}|^2 \max_{s \in (x_i, x^{\star})} |g''(s)|.$$

Recall the Taylor expansion of g around some point a:

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

What happens if we evaluate this around a fixed point x^* of g?

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

But $g(x_i) = x_{i+1}$ and $g(x^{\star}) = x^{\star}$, so

$$|x_{i+1} - x^{\star}| = |(x_i - x^{\star})g'(x^{\star}) + \frac{1}{2}(x_i - x^{\star})^2 g''(\zeta_i)|$$

$$\leq |x_i - x^{\star}||g'(x^{\star})| + \frac{1}{2}|x_i - x_{\star}|^2 \max_{s \in (x_i, x^{\star})} |g''(s)|.$$

If g has $g'(x^{\star})=0,$ we would have quadratic convergence!

Recall the Taylor expansion of g around some point a:

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

What happens if we evaluate this around a fixed point x^* of g?

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

But $g(x_i) = x_{i+1}$ and $g(x^{\star}) = x^{\star}$, so

$$|x_{i+1} - x^{\star}| = |(x_i - x^{\star})g'(x^{\star}) + \frac{1}{2}(x_i - x^{\star})^2 g''(\zeta_i)|$$

$$\leq \frac{1}{2}|x_i - x_{\star}|^2 \max_{s \in (x_i, x^{\star})} |g''(s)|.$$

If g has $g'(x^{\star})=0,$ we would have quadratic convergence!

Recall that

$$g(x) = x - \frac{f(x)}{f'(x)},$$

so (assuming $f \in C^2(\mathbb{R})$)

$$g'(x) = 1 - \left(\frac{[f'(x)]^2 - f(x)f''(x)}{[f'(x)]^2}\right)$$

Recall that

$$g(x) = x - \frac{f(x)}{f'(x)},$$

so (assuming $f \in C^2(\mathbb{R})$)

$$g'(x) = 1 - \left(\frac{[f'(x)]^2 - f(x)f''(x)}{[f'(x)]^2}\right)$$
$$= \frac{f(x)f''(x)}{[f'(x)]^2}.$$

Recall that

$$g(x) = x - \frac{f(x)}{f'(x)},$$

so (assuming $f \in C^2(\mathbb{R})$)

$$g'(x) = 1 - \left(\frac{[f'(x)]^2 - f(x)f''(x)}{[f'(x)]^2}\right)$$
$$= \frac{f(x)f''(x)}{[f'(x)]^2}.$$

If $f(x^{\star})=0$ and $f'(x^{\star})\neq 0,$ then $g'(x^{\star})=0,$ and we do get quadratic convergence!

Recall that

$$g(x) = x - \frac{f(x)}{f'(x)},$$

so (assuming $f \in C^2(\mathbb{R})$)

$$g'(x) = 1 - \left(\frac{[f'(x)]^2 - f(x)f''(x)}{[f'(x)]^2}\right)$$
$$= \frac{f(x)f''(x)}{[f'(x)]^2}.$$

If $f(x^{\star}) = 0$ and $f'(x^{\star}) \neq 0$, then $g'(x^{\star}) = 0$, and we do get quadratic convergence!

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

Take-home message

Newton's method converges quadratically to isolated roots.

Take-home message

Newton's method converges quadratically to isolated roots.

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

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

then one expects linear convergence with rate 1/2.

The true answer is $x^{\star} \approx 0.739085133215161$.

The true answer is $x^{\star} \approx 0.739085133215161$.

$$x_1 = 1 \qquad |x_1 - x^*| = 2.6 \times 10^{-1}$$

The true answer is $x^{\star} \approx 0.739085133215161$.

$$\begin{aligned} x_1 &= 1 & |x_1 - x^*| &= 2.6 \times 10^{-1} \\ x_2 &= 0.750363867840244 & |x_2 - x^*| &= 1.1 \times 10^{-2} \end{aligned}$$

The true answer is $x^* \approx 0.739085133215161$.

$$\begin{aligned} x_1 &= 1 & |x_1 - x^*| &= 2.6 \times 10^{-1} \\ x_2 &= \underline{0.750363867840244} & |x_2 - x^*| &= 1.1 \times 10^{-2} \\ x_3 &= \underline{0.739112890911362} & |x_3 - x^*| &= 2.8 \times 10^{-5} \end{aligned}$$

The true answer is $x^{\star} \approx 0.739085133215161$.

$$\begin{aligned} x_1 &= 1 & |x_1 - x^*| &= 2.6 \times 10^{-1} \\ x_2 &= \underline{0.750363867840244} & |x_2 - x^*| &= 1.1 \times 10^{-2} \\ x_3 &= \underline{0.739112890911362} & |x_3 - x^*| &= 2.8 \times 10^{-5} \\ x_4 &= \underline{0.739085133}385284 & |x_4 - x^*| &= 1.7 \times 10^{-10} \end{aligned}$$

The true answer is $x^* \approx 0.739085133215161$.

$$\begin{aligned} x_1 &= 1 & |x_1 - x^*| &= 2.6 \times 10^{-1} \\ x_2 &= \underline{0.750363867840244} & |x_2 - x^*| &= 1.1 \times 10^{-2} \\ x_3 &= \underline{0.739112890911362} & |x_3 - x^*| &= 2.8 \times 10^{-5} \\ x_4 &= \underline{0.739085133}385284 & |x_4 - x^*| &= 1.7 \times 10^{-10} \\ x_5 &= 0.739085133215161 &= x_6 &= \cdots . \end{aligned}$$

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 $x_0 \in [1/3, 1]$.

$$g(x) = x - \frac{p(x)}{p'(x)}$$

= $x - \frac{27x^3 - 27x^2 + 4}{81x^2 - 54x}$

$$g(x) = x - \frac{p(x)}{p'(x)}$$

= $x - \frac{27x^3 - 27x^2 + 4}{81x^2 - 54x}$
= $x - \frac{(3x - 2)(9x^2 - 3x - 2)}{27x(3x - 2)}$

$$g(x) = x - \frac{p(x)}{p'(x)}$$

= $x - \frac{27x^3 - 27x^2 + 4}{81x^2 - 54x}$
= $x - \frac{(3x - 2)(9x^2 - 3x - 2)}{27x(3x - 2)}$
= $\frac{2x}{3} + \frac{2}{27x} + \frac{1}{9}$.

$$g(x) = x - \frac{p(x)}{p'(x)}$$

= $x - \frac{27x^3 - 27x^2 + 4}{81x^2 - 54x}$
= $x - \frac{(3x - 2)(9x^2 - 3x - 2)}{27x(3x - 2)}$
= $\frac{2x}{3} + \frac{2}{27x} + \frac{1}{9}$.

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

$$g'(x) = \frac{2}{3} - \frac{2}{27x^2}, \quad g''(x) = \frac{4}{27x^3}.$$

$$g'(x) = \frac{2}{3} - \frac{2}{27x^2}, \quad g''(x) = \frac{4}{27x^3}.$$

We see that g''(x) > 0 on [1/3, 1] and hence g'(x) is increasing on [1/3, 1].

$$g'(x) = \frac{2}{3} - \frac{2}{27x^2}, \quad g''(x) = \frac{4}{27x^3}.$$

We see that g''(x) > 0 on [1/3, 1] and hence g'(x) is increasing on [1/3, 1]. Evaluating at the endpoints,

$$g'(1/3) = 0, \quad g'(1) = 16/27 = \gamma.$$

$$g'(x) = \frac{2}{3} - \frac{2}{27x^2}, \quad g''(x) = \frac{4}{27x^3}.$$

We see that g''(x) > 0 on [1/3, 1] and hence g'(x) is increasing on [1/3, 1]. Evaluating at the endpoints,

$$g'(1/3) = 0, \quad g'(1) = 16/27 = \gamma.$$

We now check that g is an endomorphism.

$$g'(x) = \frac{2}{3} - \frac{2}{27x^2}, \quad g''(x) = \frac{4}{27x^3}.$$

We see that g''(x) > 0 on [1/3, 1] and hence g'(x) is increasing on [1/3, 1]. Evaluating at the endpoints,

$$g'(1/3) = 0, \quad g'(1) = 16/27 = \gamma.$$

We now check that g is an endomorphism. Since $g'(x)\geq 0$ on [1/3,1], we know g is also increasing.

$$g'(x) = \frac{2}{3} - \frac{2}{27x^2}, \quad g''(x) = \frac{4}{27x^3}.$$

We see that g''(x) > 0 on [1/3, 1] and hence g'(x) is increasing on [1/3, 1]. Evaluating at the endpoints,

$$g'(1/3) = 0, \quad g'(1) = 16/27 = \gamma.$$

We now check that g is an endomorphism. Since $g'(x) \ge 0$ on [1/3, 1], we know g is also increasing.

Checking at the endpoints,

$$g(1/3) = 5/9 \in [1/3, 1], \quad g(1) = 23/27 \in [1/3, 1].$$

$$g'(x) = \frac{2}{3} - \frac{2}{27x^2}, \quad g''(x) = \frac{4}{27x^3}.$$

We see that g''(x) > 0 on [1/3, 1] and hence g'(x) is increasing on [1/3, 1]. Evaluating at the endpoints,

$$g'(1/3) = 0, \quad g'(1) = 16/27 = \gamma.$$

We now check that g is an endomorphism. Since $g'(x) \ge 0$ on [1/3, 1], we know g is also increasing.

Checking at the endpoints,

$$g(1/3) = 5/9 \in [1/3, 1], \quad g(1) = 23/27 \in [1/3, 1].$$

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

There are other fixed-point iterations for rootfinding.

Halley's method (1694)

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



Edmund Halley FRS, 1656–1742

There are other fixed-point iterations for rootfinding.

Halley's method (1694)

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



Edmund Halley FRS, 1656–1742

Halley was Savilian Professor of Geometry here in Oxford, after Wallis.
There are other fixed-point iterations for rootfinding.

Halley's method (1694)

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



Edmund Halley FRS, 1656–1742

Halley was Savilian Professor of Geometry here in Oxford, after Wallis.

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!



20 / 60

Section 2

Bonus: the secant iteration

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

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

In practice, computing derivatives of your function might be very expensive. (Think of e.g. f(x) as the evaluation of a climate model.)

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

In practice, computing derivatives of your function might be very expensive. (Think of e.g. f(x) as the evaluation of a climate model.)

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

$$x_{i+1} = g(x_i) = x - (f'(x_i))^{-1} f(x_i)$$

but we don't want to code f'(x).

$$x_{i+1} = g(x_i) = x - (f'(x_i))^{-1} f(x_i)$$

but we don't want to code f'(x).

The secant method approximates

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

with some previous data x_{i-1} .

$$x_{i+1} = g(x_i) = x - (f'(x_i))^{-1} f(x_i)$$

but we don't want to code f'(x).

The secant method approximates

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

with some previous data x_{i-1} .

This requires the user to supply both x_0 and x_{-1} .

$$x_{i+1} = g(x_i) = x - (f'(x_i))^{-1} f(x_i)$$

but we don't want to code f'(x).

The secant method approximates

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

with some previous data x_{i-1} .

This requires the user to supply both x_0 and x_{-1} .

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

$$x_{i+1} = g(x_i) = x - (f'(x_i))^{-1} f(x_i)$$

but we don't want to code f'(x).

The secant method approximates

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

with some previous data x_{i-1} .

This requires the user to supply both x_0 and x_{-1} .

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!

P. E. Farrell (Oxford)









$$\phi = \frac{1 + \sqrt{5}}{2} \approx 1.618034$$

so its convergence is superlinear, but not quite quadratic.

$$\phi = \frac{1 + \sqrt{5}}{2} \approx 1.618034$$

so its convergence is superlinear, but not quite quadratic.

The first proof to be found of this is by Terry Allen Jeeves in 1958, 300 years after Newton invented it!

$$\phi = \frac{1 + \sqrt{5}}{2} \approx 1.618034$$

so its convergence is superlinear, but not quite quadratic.

The first proof to be found of this is by Terry Allen Jeeves in 1958, 300 years after Newton invented it!

Comments on the secant method:

The method requires more information to start, and depends sensitively on it.

$$\phi = \frac{1 + \sqrt{5}}{2} \approx 1.618034$$

so its convergence is superlinear, but not quite quadratic.

The first proof to be found of this is by Terry Allen Jeeves in 1958, 300 years after Newton invented it!

Comments on the secant method:

- X The method requires more information to start, and depends sensitively on it.
- \checkmark In principle the method can be applied to nondifferentiable functions.

$$\phi = \frac{1+\sqrt{5}}{2} \approx 1.618034$$

so its convergence is superlinear, but not quite quadratic.

The first proof to be found of this is by Terry Allen Jeeves in 1958, 300 years after Newton invented it!

Comments on the secant method:

- The method requires more information to start, and depends sensitively on it.
- \checkmark 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

$$x_{i+1} = g(x_i)$$

is only converging linearly.

Suppose our fixed-point iteration

$$x_{i+1} = g(x_i)$$

is only converging linearly.

We could apply Aitken acceleration, constructing

 $x_0, x_1, x_2, x_3, x_4, \dots$ $\tilde{x}_0, \tilde{x}_1, \dots$

Suppose our fixed-point iteration

$$x_{i+1} = g(x_i)$$

is only converging linearly.

We could apply Aitken acceleration, constructing

 $x_0, x_1, x_2, x_3, x_4, \dots$ $\tilde{x}_0, \tilde{x}_1, \dots$

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

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

```
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 (x_0x_2 - x_1^2)/(x_2 - 2x_1 + x_0)

end while

return g(x)

end function
```

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

```
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 (x_0x_2 - x_1^2)/(x_2 - 2x_1 + x_0)

end while

return g(x)

end function
```

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

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

```
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 (x_0x_2 - x_1^2)/(x_2 - 2x_1 + x_0)

end while

return g(x)

end function
```

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

Does this really help?

| P. E. Farrell (| Oxford |
|-----------------|--------|
|-----------------|--------|

Yes, it does, under certain conditions:

Steffensen's theorem (1933)

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

This can achieve quadratic convergence, without derivatives!

Let's see two examples.

Let's see two examples.

We previously considered the fixed-point iteration

$$g(x) = \frac{x+1}{x}$$

for calculating the golden ratio ϕ .

Let's see two examples.

We previously considered the fixed-point iteration

$$g(x) = \frac{x+1}{x}$$

for calculating the golden ratio ϕ .

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

$$f(x) = (x-1)^2.$$

Let's apply Newton's method to

$$f(x) = (x-1)^2.$$

This gives

$$g(x) = x - \frac{(x-1)^2}{2x-2}.$$

Since f'(1) = 0, $g'(1) \neq 0$, and we only achieve linear convergence:

Let's apply Newton's method to

$$f(x) = (x-1)^2.$$

This gives

$$g(x) = x - \frac{(x-1)^2}{2x-2}.$$

Since f'(1) = 0, $g'(1) \neq 0$, and we only achieve linear convergence:

```
In [17]: newton(lambda x: (x-1)**2, lambda x: 2*x - 2, 0, 1e-4)

Iteration 0: x = 0.000000e+00 f(x) = 1.000000e+00

Iteration 1: x = 5.000000e-01 f(x) = 2.500000e-01

Iteration 2: x = 7.50000e-01 f(x) = 6.250000e-02

Iteration 3: x = 8.75000e-01 f(x) = 1.562500e-02

Iteration 4: x = 9.375000e-01 f(x) = 3.906250e-03

Iteration 5: x = 9.687500e-01 f(x) = 9.765625e-04

Iteration 6: x = 9.843750e-01 f(x) = 2.441406e-04

Iteration 7: x = 9.921875e-01 f(x) = 6.103516e-05

Out[17]: 0.9921875
```

P. E. Farrell (Oxford)

Converging linearly, you say?

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.

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

Philosophical remark

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

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

Philosophical remark

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

1. an efficient evaluation strategy for polynomials in the monomial basis;

- 1. an efficient evaluation strategy for polynomials in the monomial basis;
- 2. an iteration scheme for finding the roots of polynomials that combines Newton's method with the evaluation scheme.

- 1. an efficient evaluation strategy for polynomials in the monomial basis;
- 2. an iteration scheme for finding the roots of polynomials that combines Newton's method with the evaluation scheme.

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.

- 1. an efficient evaluation strategy for polynomials in the monomial basis;
- 2. an iteration scheme for finding the roots of polynomials that combines Newton's method with the evaluation scheme.

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.

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

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

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

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

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

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

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

multiplications. Scaling like n^2 is bad!

 $a_0 + a_1 x + \dots + a_n x^n$

$$a_0 + a_1 x + \dots + a_n x^n = a_0 + \dots$$

$$a_0 + a_1 x + \dots + a_n x^n = a_0 + x$$
 (

$$a_0 + a_1 x + \dots + a_n x^n = a_0 + x (a_1 + x)$$

$$a_0 + a_1 x + \dots + a_n x^n = a_0 + x (a_1 + x (a_2 + \dots + x (a_{n-1} + xa_n) \dots)).$$

$$a_0 + a_1 x + \dots + a_n x^n = a_0 + x (a_1 + x (a_2 + \dots + x (a_{n-1} + xa_n) \dots))$$

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

$$a_0 + a_1 x + \dots + a_n x^n = a_0 + x (a_1 + x (a_2 + \dots + x (a_{n-1} + xa_n) \dots)).$$

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

$$b_n \coloneqq a_n$$

$$a_0 + a_1 x + \dots + a_n x^n = a_0 + x (a_1 + x (a_2 + \dots + x (a_{n-1} + xa_n) \dots))$$

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

$$b_n \coloneqq a_n$$
$$b_{n-1} \coloneqq a_{n-1} + b_n r$$

$$a_0 + a_1 x + \dots + a_n x^n = a_0 + x (a_1 + x (a_2 + \dots + x (a_{n-1} + xa_n) \dots))$$

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

$$b_n \coloneqq a_n$$
$$b_{n-1} \coloneqq a_{n-1} + b_n r$$
$$\vdots$$
$$b_i \coloneqq a_i + b_{i+1} r$$

$$a_0 + a_1 x + \dots + a_n x^n = a_0 + x (a_1 + x (a_2 + \dots + x (a_{n-1} + xa_n) \dots))$$

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

$$b_n \coloneqq a_n$$

$$b_{n-1} \coloneqq a_{n-1} + b_n r$$

$$\vdots$$

$$b_i \coloneqq a_i + b_{i+1} r$$

$$\vdots$$

$$b_1 \coloneqq a_1 + b_2 r$$

$$a_0 + a_1 x + \dots + a_n x^n = a_0 + x (a_1 + x (a_2 + \dots + x (a_{n-1} + xa_n) \dots)).$$

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

Algorithmically, to evaluate p(r) for given $r\in\mathbb{R}$ we calculate

$$b_n \coloneqq a_n$$

$$b_{n-1} \coloneqq a_{n-1} + b_n r$$

$$\vdots$$

$$b_i \coloneqq a_i + b_{i+1} r$$

$$\vdots$$

$$b_1 \coloneqq a_1 + b_2 r$$

$$b_0 \coloneqq a_0 + b_1 r.$$

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

Theorem

Define the polynomial

$$Q(x) \coloneqq b_n x^{n-1} + b_{n-1} x^{n-2} + \dots + b_2 x + b_1.$$

Then

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

Theorem

Define the polynomial

$$Q(x) \coloneqq b_n x^{n-1} + b_{n-1} x^{n-2} + \dots + b_2 x + b_1.$$

Then

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

Before proving this, note that indeed $p(r) = b_0$, and

$$p'(x) = Q(x) + (x - r)Q'(x),$$

Theorem

Define the polynomial

$$Q(x) \coloneqq b_n x^{n-1} + b_{n-1} x^{n-2} + \dots + b_2 x + b_1.$$

Then

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

Before proving this, note that indeed $p(r) = b_0$, and

$$p'(x) = Q(x) + (x - r)Q'(x),$$

so in particular

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

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

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

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

$$(x-r)Q(x) + b_0 = (x-r)(b_n x^{n-1} + \dots + b_1) + b_0$$

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

$$(x-r)Q(x) + b_0 = (x-r)(b_n x^{n-1} + \dots + b_1) + b_0$$

= $x(b_n x^{n-1} + \dots + b_1) - r(b_n x^{n-1} + \dots + b_1) + b_0$

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

$$(x-r)Q(x) + b_0 = (x-r)(b_n x^{n-1} + \dots + b_1) + b_0$$

= $x(b_n x^{n-1} + \dots + b_1) - r(b_n x^{n-1} + \dots + b_1) + b_0$
= $b_n x^n + (b_{n-1} - b_n r) x^{n-1} + \dots + (b_0 - b_1 r)$

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

$$(x-r)Q(x) + b_0 = (x-r)(b_n x^{n-1} + \dots + b_1) + b_0$$

= $x(b_n x^{n-1} + \dots + b_1) - r(b_n x^{n-1} + \dots + b_1) + b_0$
= $b_n x^n + (b_{n-1} - b_n r)x^{n-1} + \dots + (b_0 - b_1 r)$
= $a_n x^n + a_{n-1} x^{n-1} + \dots + a_0$,

since $a_i = b_i - b_{i+1}r$ for i < n.

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

$$(x-r)Q(x) + b_0 = (x-r)(b_n x^{n-1} + \dots + b_1) + b_0$$

= $x(b_n x^{n-1} + \dots + b_1) - r(b_n x^{n-1} + \dots + b_1) + b_0$
= $b_n x^n + (b_{n-1} - b_n r)x^{n-1} + \dots + (b_0 - b_1 r)$
= $a_n x^n + a_{n-1} x^{n-1} + \dots + a_0$,

since $a_i = b_i - b_{i+1}r$ for i < n.

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

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

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

$$(x-r)Q(x) + b_0 = (x-r)(b_n x^{n-1} + \dots + b_1) + b_0$$

= $x(b_n x^{n-1} + \dots + b_1) - r(b_n x^{n-1} + \dots + b_1) + b_0$
= $b_n x^n + (b_{n-1} - b_n r)x^{n-1} + \dots + (b_0 - b_1 r)$
= $a_n x^n + a_{n-1} x^{n-1} + \dots + a_0$,

since $a_i = b_i - b_{i+1}r$ for i < n.

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

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

function horner($[a_0, \cdots, a_n]$, x_0 , tol, maxit) $x \leftarrow x_0$

function horner($[a_0, \cdots, a_n]$, x_0 , tol, maxit) $x \leftarrow x_0$ for $i = 1, \dots, \text{maxit}$ do $b \leftarrow a_n x + a_{n-1}$ # He

Horner eval for p
function horner($[a_0, \dots, a_n]$, x_0 , tol, maxit) $x \leftarrow x_0$ for $i = 1, \dots$, maxit do $b \leftarrow a_n x + a_{n-1}$ $c \leftarrow a_n$ for $k = n - 1, n - 2, \dots, 1, 0$ do $c \leftarrow cx + b$

 $\begin{array}{l} \# \mbox{ Horner eval for } p \\ \# \mbox{ Horner eval for } p' \end{array}$

function horner($[a_0, \dots, a_n]$, x_0 , tol, maxit) $x \leftarrow x_0$ for $i = 1, \dots$, maxit do $b \leftarrow a_n x + a_{n-1}$ $c \leftarrow a_n$ for $k = n - 1, n - 2, \dots, 1, 0$ do $c \leftarrow cx + b$ $b \leftarrow bx + a_i$ end for

 $\begin{array}{l} \# \mbox{ Horner eval for } p \\ \# \mbox{ Horner eval for } p' \end{array}$

function horner($[a_0, \cdots, a_n]$, x_0 , tol, maxit) $x \leftarrow x_0$ for $i = 1, \ldots, \text{maxit} do$ $b \leftarrow a_n x + a_{n-1}$ # Horner eval for p# Horner eval for p' $c \leftarrow a_n$ for $k = n - 1, n - 2, \dots, 1, 0$ do $c \leftarrow cx + b$ $b \leftarrow bx + a_i$ end for if |b| < tol then# success return x end if

function horner($[a_0, \cdots, a_n]$, x_0 , tol, maxit) $x \leftarrow x_0$ for $i = 1, \ldots, \text{maxit}$ do # Horner eval for p $b \leftarrow a_n x + a_{n-1}$ # Horner eval for p' $c \leftarrow a_n$ for $k = n - 1, n - 2, \dots, 1, 0$ do $c \leftarrow cx + b$ $b \leftarrow bx + a_i$ end for 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

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

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

 $|f(n)| \leq Mg(n)$ for all $n \geq n_0$.

We can summarise with the following useful notation:

Definition (Big \mathcal{O} notation)

For g(n) > 0, we say

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

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

$$|f(n)| \leq Mg(n)$$
 for all $n \geq n_0$.

The number of operations to evaluate a degree-n polynomial is: • $\mathcal{O}(n^2)$ for the naïve way, but We can summarise with the following useful notation:

Definition (Big \mathcal{O} notation)

For g(n) > 0, we say

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

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

$$|f(n)| \leq Mg(n)$$
 for all $n \geq n_0$.

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

Any algorithm for evaluating an arbitrary polynomial must require at least n additions (Ostrowski, 1954) and at least n multiplications (Pan, 1966).

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

Theorem

Any algorithm for evaluating an arbitrary polynomial must require at least n additions (Ostrowski, 1954) and at least n multiplications (Pan, 1966).

Since Horner's scheme employs n additions and n multiplications, it is optimal (for arbitrary polynomials).

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

Theorem

Any algorithm for evaluating an arbitrary polynomial must require at least n additions (Ostrowski, 1954) and at least n multiplications (Pan, 1966).

Since Horner's scheme employs n additions and n multiplications, it is optimal (for arbitrary polynomials).

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

Philosophical remark

Equivalent expressions can have different algorithmic properties!

Philosophical remark

Equivalent expressions can have different algorithmic properties!

In Horner's case, we had

 $a_0 + a_1 x + \dots + a_n x^n = a_0 + x (a_1 + x (a_2 + \dots + x (a_{n-1} + xa_n) \dots)).$

Philosophical remark

Equivalent expressions can have different algorithmic properties!

In Horner's case, we had

 $a_0 + a_1 x + \dots + a_n x^n = a_0 + x (a_1 + x (a_2 + \dots + x (a_{n-1} + xa_n) \dots)).$

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

- Does the algorithm terminate?
- Does the algorithm give the correct answer?
- How fast does the algorithm converge to the answer?
- How many operations does it take?

- Does the algorithm terminate?
- Does the algorithm give the correct answer?
- How fast does the algorithm converge to the answer?
- How many operations does it take?

There's another very important question we might want to ask:

- Does the algorithm terminate?
- Does the algorithm give the correct answer?
- How fast does the algorithm converge to the answer?
- How many operations does it take?

There's another very important question we might want to ask:

Can we parallelise the algorithm?

- Does the algorithm terminate?
- Does the algorithm give the correct answer?
- How fast does the algorithm converge to the answer?
- How many operations does it take?

There's another very important question we might want to ask:

Can we parallelise the algorithm?

Every computer nowadays has multiple processing units. (My phone has 8.) Can we use them?

$$a_0 + a_1 x + \dots + a_n x^n$$

= $(a_0 + a_2 x^2 + a_4 x^4 + \dots) + (a_1 x + a_3 x^3 + a_5 x^5 + \dots)$

$$a_0 + a_1 x + \dots + a_n x^n$$

= $(a_0 + a_2 x^2 + a_4 x^4 + \dots) + (a_1 x + a_3 x^3 + a_5 x^5 + \dots)$
= $(a_0 + a_2 x^2 + a_4 x^4 + \dots) + x (a_1 + a_3 x^2 + a_5 x^4 + \dots)$

$$a_0 + a_1 x + \dots + a_n x^n$$

= $(a_0 + a_2 x^2 + a_4 x^4 + \dots) + (a_1 x + a_3 x^3 + a_5 x^5 + \dots)$
= $(a_0 + a_2 x^2 + a_4 x^4 + \dots) + x (a_1 + a_3 x^2 + a_5 x^4 + \dots)$
= $p_1(x^2) + x p_2(x^2)$

$$a_0 + a_1 x + \dots + a_n x^n$$

= $(a_0 + a_2 x^2 + a_4 x^4 + \dots) + (a_1 x + a_3 x^3 + a_5 x^5 + \dots)$
= $(a_0 + a_2 x^2 + a_4 x^4 + \dots) + x (a_1 + a_3 x^2 + a_5 x^4 + \dots)$
= $p_1(x^2) + x p_2(x^2)$

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

$$a_0 + a_1 x + \dots + a_n x^n$$

= $(a_0 + a_2 x^2 + a_4 x^4 + \dots) + (a_1 x + a_3 x^3 + a_5 x^5 + \dots)$
= $(a_0 + a_2 x^2 + a_4 x^4 + \dots) + x (a_1 + a_3 x^2 + a_5 x^4 + \dots)$
= $p_1(x^2) + x p_2(x^2)$

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

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

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.

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

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

$$p_1(x) = \frac{p_0(x)}{(x - x^*)}$$

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.

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

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

$$p_1(x) = \frac{p_0(x)}{(x - x^*)}$$

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.

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 diagonalisable matrices:

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

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

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

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-

Vera Kublanovskaya, 1920–2012

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

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

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 $20^{\rm th}$ century.



John Francis, 1934-

Vera Kublanovskaya, 1920–2012

You can learn more in A7: Numerical Analysis.

P. E. Farrell (Oxford)

Newton's method

Given p(x) of degree n, we want to construct an A with characteristic polynomial p(x).

Given p(x) of degree n, we want to construct an A with characteristic polynomial p(x). Let

$$p(x) = a_0 + a_1 x + \dots + x^n$$

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

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

.

Given p(x) of degree n, we want to construct an A with characteristic polynomial p(x). Let

$$p(x) = a_0 + a_1 x + \dots + x^n$$

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

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

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

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

By applying the QR algorithm for eigenvalues to the companion matrix, we can find all roots in $\mathcal{O}(n^3)$ operations.

By applying the QR algorithm for eigenvalues to the companion matrix, we can find all roots in $\mathcal{O}(n^3)$ operations.

We previously saw that Newton's method can get stuck in a cycle for $p(x) = x^3 - 2x + 2$. No problem:
Section 8

Bonus: representing polynomials

Algorithms are usually tied to the *data structures* we use.

Algorithms are usually tied to the *data structures* we use.

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 $\{M_i\}$:

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

Algorithms are usually tied to the *data structures* we use.

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 $\{M_i\}$:

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

The algorithms take in $[a_0, a_1, \ldots, a_n] \in \mathbb{R}^{n+1}$ to represent p.

Algorithms are usually tied to the *data structures* we use.

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 $\{M_i\}$:

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

The algorithms take in $[a_0, a_1, \ldots, a_n] \in \mathbb{R}^{n+1}$ to represent p.

A natural question to ask:

is the map $a \mapsto p$ stable?

If we make a perturbation δa to a, how big can the perturbation δp be? For the monomial basis $\{M_i\}$, the answer is very very big: If we make a perturbation δa to a, how big can the perturbation δp be? For the monomial basis $\{M_i\}$, the answer is *very very big*:

Construct

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

then perturb its monomial coefficients by

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



James H. Wilkinson, 1919-1986

If we make a perturbation δa to a, how big can the perturbation δp be? For the monomial basis $\{M_i\}$, the answer is *very very big*:

Construct

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

then perturb its monomial coefficients by

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



James H. Wilkinson, 1919-1986

The resulting δp has

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

for a stability constant of

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

Not all bases are equally good.

Not all bases are equally good.

For example, for $\varepsilon > 0$, the set

$$\{(1,0)^{\top},(1,\varepsilon)^{\top}\}$$

is as much of a basis for \mathbb{R}^2 as

 $\{(1,0)^{\top}, (0,1)^{\top}\}.$

Not all bases are equally good.

For example, for $\varepsilon > 0$, the set

$$\{(1,0)^{\top},(1,\varepsilon)^{\top}\}$$

is as much of a basis for \mathbb{R}^2 as

 $\{(1,0)^{\top}, (0,1)^{\top}\}.$

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 $\{T_i: [-1,1] \rightarrow [-1,1]\}$ satisfy

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

The role of the \hat{x} is to map the input interval [a, b] to [-1, 1].

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 $\{T_i: [-1,1] \rightarrow [-1,1]\}$ satisfy

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

The role of the \hat{x} is to map the input interval [a, b] to [-1, 1].



Chebyshev polynomials. Credit: Glosser.ca, Wikipedia

P. E. Farrell (Oxford)

Just as a polynomial p has a finite Chebyshev series, general functions f have infinite Chebyshev series. These expansions converge very, very fast:

Just as a polynomial p has a finite Chebyshev series, general functions f have infinite Chebyshev series. These expansions converge very, very fast:

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

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

Just as a polynomial p has a finite Chebyshev series, general functions f have infinite Chebyshev series. These expansions converge very, very fast:

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

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

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

But we can't apply Horner's method for evaluation, or the companion matrix trick for rootfinding, to Chebyshev expansions!

But we can't apply Horner's method for evaluation, or the companion matrix trick for rootfinding, to Chebyshev expansions!

Our algorithms are tied to our choice of representation: to our choice of basis, or (from a CS perspective) the data structure we use.

But we can't apply Horner's method for evaluation, or the companion matrix trick for rootfinding, to Chebyshev expansions!

Our algorithms are tied to our choice of representation: to our choice of basis, or (from a CS perspective) the data structure we use.

Good news

For Chebyshev bases, analogous algorithms exist:

- ✓ the second barycentric formula, for 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} :

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

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

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

- 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

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

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

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

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.

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.

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

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

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.

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

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

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

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

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

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

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

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

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

•

 $F(\mathbf{x}) = F(\mathbf{a}) + DF(\mathbf{a})(\mathbf{x} - \mathbf{a}) + \text{ higher-order terms.}$

$$F(\mathbf{x}) = F(\mathbf{a}) + DF(\mathbf{a})(\mathbf{x} - \mathbf{a}) + \text{ higher-order terms.}$$

Following the reasoning from one dimension,

$$F(\mathbf{x}_{i+1}) = F(\mathbf{x}_i + \delta \mathbf{x}) \approx F(\mathbf{x}_i) + DF(\mathbf{x}_i)\delta \mathbf{x},$$

$$F(\mathbf{x}) = F(\mathbf{a}) + DF(\mathbf{a})(\mathbf{x} - \mathbf{a}) + \text{ higher-order terms.}$$

Following the reasoning from one dimension,

$$F(\mathbf{x}_{i+1}) = F(\mathbf{x}_i + \delta \mathbf{x}) \approx F(\mathbf{x}_i) + DF(\mathbf{x}_i)\delta \mathbf{x},$$

and optimistically setting $F(\mathbf{x}_{i+1}) = 0$, we get

$$\delta \mathbf{x} = -[DF(\mathbf{x}_i)]^{-1}F(\mathbf{x}_i), \quad \mathbf{x}_{i+1} = \mathbf{x}_i + \delta \mathbf{x}.$$

$$F(\mathbf{x}) = F(\mathbf{a}) + DF(\mathbf{a})(\mathbf{x} - \mathbf{a}) + \text{ higher-order terms.}$$

Following the reasoning from one dimension,

$$F(\mathbf{x}_{i+1}) = F(\mathbf{x}_i + \delta \mathbf{x}) \approx F(\mathbf{x}_i) + DF(\mathbf{x}_i)\delta \mathbf{x},$$

and optimistically setting $F(\mathbf{x}_{i+1}) = 0$, we get

$$\delta \mathbf{x} = -[DF(\mathbf{x}_i)]^{-1}F(\mathbf{x}_i), \quad \mathbf{x}_{i+1} = \mathbf{x}_i + \delta \mathbf{x}.$$

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

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

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

P. E. Farrell (Oxford)

Newton-Raphson method

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

Comments:

✓ Still a fixed-point method.
$$\mathbf{x}_{i+1} = g(\mathbf{x}_i) \coloneqq \mathbf{x}_i - (DF(\mathbf{x}_i))^{-1} F(\mathbf{x}_i).$$

- ✓ Still a fixed-point method.
- ✓ Fixed points of g are roots of F.

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

- ✓ Still a fixed-point method.
- ✓ Fixed points of g are roots of F.
- \checkmark We still require F to be differentiable.

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

- ✓ Still a fixed-point method.
- \checkmark Fixed points of g are roots of F.
- \checkmark We still require F to be differentiable.
- \checkmark We now require DF to be invertible at every iterate.

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

- ✓ Still a fixed-point method.
- \checkmark Fixed points of g are roots of F.
- \checkmark We still require F to be differentiable.
- \checkmark We now require DF to be invertible at every iterate.
- × We have to solve linear systems (worse case $\mathcal{O}(N^3)$ operations).

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

- ✓ Still a fixed-point method.
- ✓ Fixed points of g are roots of F.
- \checkmark We still require F to be differentiable.
- \checkmark We now require DF to be invertible at every iterate.
- × We have to solve linear systems (worse case $\mathcal{O}(N^3)$ operations).
- ✓ Sometimes the linear systems can be solved in $\mathcal{O}(N)$ operations.

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

- ✓ Still a fixed-point method.
- ✓ Fixed points of g are roots of F.
- \mathbf{X} We still require F to be differentiable.
- \checkmark We now require DF to be invertible at every iterate.
- × We have to solve linear systems (worse case $\mathcal{O}(N^3)$ operations).
- ✓ Sometimes the linear systems can be solved in $\mathcal{O}(N)$ operations.
- ✓ If x_0 is close to x^* , Newton's method usually converges quadratically.

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

- ✓ Still a fixed-point method.
- ✓ Fixed points of g are roots of F.
- \mathbf{X} We still require F to be differentiable.
- \checkmark We now require DF to be invertible at every iterate.
- × We have to solve linear systems (worse case $\mathcal{O}(N^3)$ operations).
- ✓ Sometimes the linear systems can be solved in O(N) operations.
- ✓ If x_0 is close to x^* , Newton's method usually converges quadratically.
- \checkmark If x_0 is far away, the method can diverge or get stuck in a cycle.

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

- ✓ Still a fixed-point method.
- ✓ Fixed points of g are roots of F.
- \mathbf{X} We still require F to be differentiable.
- \checkmark We now require DF to be invertible at every iterate.
- × We have to solve linear systems (worse case $\mathcal{O}(N^3)$ operations).
- ✓ Sometimes the linear systems can be solved in $\mathcal{O}(N)$ operations.
- ✓ If x_0 is close to x^* , Newton's method usually converges quadratically.
- \checkmark 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

Let's do an example. Take

$$F(x,y) = \begin{pmatrix} xy + y^2 - 2 \\ x^3y - 3x - 1 \end{pmatrix},$$

Let's do an example. Take

$$F(x,y) = \begin{pmatrix} xy + y^2 - 2 \\ x^3y - 3x - 1 \end{pmatrix}, \text{ with } DF(x,y) = \begin{pmatrix} y & x + 2y \\ 3x^2y - 3 & x^3 \end{pmatrix}.$$

Let's do an example. Take

$$F(x,y) = \begin{pmatrix} xy + y^2 - 2 \\ x^3y - 3x - 1 \end{pmatrix}, \text{ with } DF(x,y) = \begin{pmatrix} y & x + 2y \\ 3x^2y - 3 & x^3 \end{pmatrix}.$$

Starting at $\mathbf{x}_0 = (0,1)^\top$, we have to solve

$$\begin{pmatrix} 1 & 2 \\ -3 & 0 \end{pmatrix} \begin{pmatrix} \delta x \\ \delta y \end{pmatrix} = \begin{pmatrix} 1 \\ 1 \end{pmatrix}.$$

Let's do an example. Take

$$F(x,y) = \begin{pmatrix} xy + y^2 - 2 \\ x^3y - 3x - 1 \end{pmatrix}, \text{ with } DF(x,y) = \begin{pmatrix} y & x + 2y \\ 3x^2y - 3 & x^3 \end{pmatrix}.$$

Starting at $\mathbf{x}_0 = (0,1)^\top$, we have to solve

$$\begin{pmatrix} 1 & 2 \\ -3 & 0 \end{pmatrix} \begin{pmatrix} \delta x \\ \delta y \end{pmatrix} = \begin{pmatrix} 1 \\ 1 \end{pmatrix}.$$

This yields $(\delta x, \delta y)^\top = (-1/3, 2/3)^\top$, so

$$\mathbf{x}_1 = \mathbf{x}_0 + \delta \mathbf{x} = \begin{pmatrix} 0\\1 \end{pmatrix} + \begin{pmatrix} -1/3\\2/3 \end{pmatrix} = \begin{pmatrix} -1/3\\5/3 \end{pmatrix}$$

Repeating the procedure, the next iterates are

$$\mathbf{x}_2 = \begin{pmatrix} -0.357668\\ 1.606112 \end{pmatrix},$$

Let's do an example. Take

$$F(x,y) = \begin{pmatrix} xy + y^2 - 2 \\ x^3y - 3x - 1 \end{pmatrix}, \text{ with } DF(x,y) = \begin{pmatrix} y & x + 2y \\ 3x^2y - 3 & x^3 \end{pmatrix}.$$

Starting at $\mathbf{x}_0 = (0,1)^\top$, we have to solve

$$\begin{pmatrix} 1 & 2 \\ -3 & 0 \end{pmatrix} \begin{pmatrix} \delta x \\ \delta y \end{pmatrix} = \begin{pmatrix} 1 \\ 1 \end{pmatrix}.$$

This yields $(\delta x,\delta y)^{\top}=(-1/3,2/3)^{\top}$, so

$$\mathbf{x}_1 = \mathbf{x}_0 + \delta \mathbf{x} = \begin{pmatrix} 0\\1 \end{pmatrix} + \begin{pmatrix} -1/3\\2/3 \end{pmatrix} = \begin{pmatrix} -1/3\\5/3 \end{pmatrix}.$$

Repeating the procedure, the next iterates are

$$\mathbf{x}_2 = \begin{pmatrix} -0.357668\\ 1.606112 \end{pmatrix}, \quad \mathbf{x}_3 = \begin{pmatrix} -0.357838\\ 1.604407 \end{pmatrix},$$

Let's do an example. Take

$$F(x,y) = \begin{pmatrix} xy + y^2 - 2 \\ x^3y - 3x - 1 \end{pmatrix}, \text{ with } DF(x,y) = \begin{pmatrix} y & x + 2y \\ 3x^2y - 3 & x^3 \end{pmatrix}.$$

Starting at $\mathbf{x}_0 = (0,1)^\top$, we have to solve

$$\begin{pmatrix} 1 & 2 \\ -3 & 0 \end{pmatrix} \begin{pmatrix} \delta x \\ \delta y \end{pmatrix} = \begin{pmatrix} 1 \\ 1 \end{pmatrix}.$$

This yields $(\delta x,\delta y)^{\top}=(-1/3,2/3)^{\top}$, so

$$\mathbf{x}_1 = \mathbf{x}_0 + \delta \mathbf{x} = \begin{pmatrix} 0\\1 \end{pmatrix} + \begin{pmatrix} -1/3\\2/3 \end{pmatrix} = \begin{pmatrix} -1/3\\5/3 \end{pmatrix}.$$

Repeating the procedure, the next iterates are

$$\mathbf{x}_2 = \begin{pmatrix} -0.357668\\ 1.606112 \end{pmatrix}, \quad \mathbf{x}_3 = \begin{pmatrix} -0.357838\\ 1.604407 \end{pmatrix}, \quad \mathbf{x}_4 = \begin{pmatrix} -0.357838\\ 1.604406 \end{pmatrix}$$

.

Section 4

Convergence

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

Given $\mathbf{x} \in \mathbb{R}^N$, we define its ∞ -norm to be

$$\|\mathbf{x}\|_{\infty} \coloneqq \max_{k=1,\dots,N} |\mathbf{x}^k|.$$

Convergence

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

Given $\mathbf{x} \in \mathbb{R}^N$, we define its ∞ -norm to be

$$\|\mathbf{x}\|_{\infty} \coloneqq \max_{k=1,\dots,N} |\mathbf{x}^k|.$$

Definition (Convergence of a vector-valued sequence)

We say
$$(\mathbf{x}_i) \to \mathbf{x}^*$$
 in the $\|\cdot\|_{\infty}$ norm if
$$\lim_{i \to \infty} \|\mathbf{x}_i - \mathbf{x}^*\|_{\infty} = 0.$$

Convergence

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

Given $\mathbf{x} \in \mathbb{R}^N$, we define its ∞ -norm to be

$$\|\mathbf{x}\|_{\infty} \coloneqq \max_{k=1,\dots,N} |\mathbf{x}^k|.$$

Definition (Convergence of a vector-valued sequence)

We say
$$(\mathbf{x}_i) \to \mathbf{x}^*$$
 in the $\|\cdot\|_{\infty}$ norm if
$$\lim_{i \to \infty} \|\mathbf{x}_i - \mathbf{x}^*\|_{\infty} = 0.$$

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).

Assuming Newton's method converges, how fast does it converge? From our one-dimensional experience, we expect quadratic convergence to isolated roots. Assuming Newton's method converges, how fast does it converge? From our one-dimensional experience, we expect quadratic convergence to isolated roots.

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 $\mathbf{x}^* \in \mathbb{R}^N$ is an isolated root of F, i.e. $F(\mathbf{x}^*) = \mathbf{0}$ with $DF(\mathbf{x}^*)$ nonsingular. Then if \mathbf{x}_0 is close enough to \mathbf{x}^* , the Newton sequence will converge quadratically.

Assuming Newton's method converges, how fast does it converge? From our one-dimensional experience, we expect quadratic convergence to isolated roots.

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 $\mathbf{x}^* \in \mathbb{R}^N$ is an isolated root of F, i.e. $F(\mathbf{x}^*) = \mathbf{0}$ with $DF(\mathbf{x}^*)$ nonsingular. Then if \mathbf{x}_0 is close enough to \mathbf{x}^* , the Newton sequence will converge quadratically.

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

Section 5

Bonus: Affine covariance

Given $F: \mathbb{R}^N \to \mathbb{R}^N$, and $\mathbf{x}_0 \in \mathbb{R}^N$, we construct the sequence $\mathbf{x}_0, \mathbf{x}_1, \ldots$.

Given $F: \mathbb{R}^N \to \mathbb{R}^N$, and $\mathbf{x}_0 \in \mathbb{R}^N$, we construct the sequence $\mathbf{x}_0, \mathbf{x}_1, \dots$.

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}^* .

Given $F: \mathbb{R}^N \to \mathbb{R}^N$, and $\mathbf{x}_0 \in \mathbb{R}^N$, we construct the sequence $\mathbf{x}_0, \mathbf{x}_1, \dots$.

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}^* .

Theorem (Affine covariance)

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

$$\tilde{\mathbf{x}}_0, \tilde{\mathbf{x}}_1, \tilde{\mathbf{x}}_2, \dots$$

Proof.

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

Assume $\mathbf{x}_i = \tilde{\mathbf{x}}_i$ at iteration i. Then the Newton update for \tilde{F} satisfies $-\delta \tilde{\mathbf{x}}_i = [D\tilde{F}(\tilde{\mathbf{x}}_i)]^{-1}\tilde{F}(\tilde{\mathbf{x}}_i) = [ADF(\mathbf{x}_i)]^{-1}AF(\mathbf{x}_i)$

$$\tilde{\mathbf{x}}_0, \tilde{\mathbf{x}}_1, \tilde{\mathbf{x}}_2, \dots$$

Proof.

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

Assume $\mathbf{x}_i = \tilde{\mathbf{x}}_i$ at iteration i. Then the Newton update for \tilde{F} satisfies $-\delta \tilde{\mathbf{x}}_i = [D\tilde{F}(\tilde{\mathbf{x}}_i)]^{-1}\tilde{F}(\tilde{\mathbf{x}}_i) = [ADF(\mathbf{x}_i)]^{-1}AF(\mathbf{x}_i)$ $= [DF(\mathbf{x}_i)]^{-1}A^{-1}AF(\mathbf{x}_i)$

$$\tilde{\mathbf{x}}_0, \tilde{\mathbf{x}}_1, \tilde{\mathbf{x}}_2, \dots$$

Proof.

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

Assume $\mathbf{x}_i = \tilde{\mathbf{x}}_i$ at iteration i. Then the Newton update for \tilde{F} satisfies $-\delta \tilde{\mathbf{x}}_i = [D\tilde{F}(\tilde{\mathbf{x}}_i)]^{-1}\tilde{F}(\tilde{\mathbf{x}}_i) = [ADF(\mathbf{x}_i)]^{-1}AF(\mathbf{x}_i)$ $= [DF(\mathbf{x}_i)]^{-1}A^{-1}AF(\mathbf{x}_i)$ $= [DF(\mathbf{x}_i)]^{-1}F(\mathbf{x}_i)$

$$\tilde{\mathbf{x}}_0, \tilde{\mathbf{x}}_1, \tilde{\mathbf{x}}_2, \dots$$

Proof.

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

Assume $\mathbf{x}_i = \tilde{\mathbf{x}}_i$ at iteration i. Then the Newton update for \tilde{F} satisfies $-\delta \tilde{\mathbf{x}}_i = [D\tilde{F}(\tilde{\mathbf{x}}_i)]^{-1}\tilde{F}(\tilde{\mathbf{x}}_i) = [ADF(\mathbf{x}_i)]^{-1}AF(\mathbf{x}_i)$ $= [DF(\mathbf{x}_i)]^{-1}A^{-1}AF(\mathbf{x}_i)$ $= [DF(\mathbf{x}_i)]^{-1}F(\mathbf{x}_i) = -\delta \mathbf{x}_i.$

$$\tilde{\mathbf{x}}_0, \tilde{\mathbf{x}}_1, \tilde{\mathbf{x}}_2, \dots$$

Proof.

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

Assume $\mathbf{x}_i = \tilde{\mathbf{x}}_i$ at iteration i. Then the Newton update for \tilde{F} satisfies $-\delta \tilde{\mathbf{x}}_i = [D\tilde{F}(\tilde{\mathbf{x}}_i)]^{-1}\tilde{F}(\tilde{\mathbf{x}}_i) = [ADF(\mathbf{x}_i)]^{-1}AF(\mathbf{x}_i)$ $= [DF(\mathbf{x}_i)]^{-1}A^{-1}AF(\mathbf{x}_i)$ $= [DF(\mathbf{x}_i)]^{-1}F(\mathbf{x}_i) = -\delta \mathbf{x}_i.$

Hence $\mathbf{x}_{i+1} = \tilde{\mathbf{x}}_{i+1}$, and the result follows by induction.

$$\tilde{\mathbf{x}}_0, \tilde{\mathbf{x}}_1, \tilde{\mathbf{x}}_2, \dots$$

Proof.

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

Assume $\mathbf{x}_i = \tilde{\mathbf{x}}_i$ at iteration i. Then the Newton update for \tilde{F} satisfies $-\delta \tilde{\mathbf{x}}_i = [D\tilde{F}(\tilde{\mathbf{x}}_i)]^{-1}\tilde{F}(\tilde{\mathbf{x}}_i) = [ADF(\mathbf{x}_i)]^{-1}AF(\mathbf{x}_i)$ $= [DF(\mathbf{x}_i)]^{-1}A^{-1}AF(\mathbf{x}_i)$ $= [DF(\mathbf{x}_i)]^{-1}F(\mathbf{x}_i) = -\delta \mathbf{x}_i.$

Hence $\mathbf{x}_{i+1} = \tilde{\mathbf{x}}_{i+1}$, and the result follows by induction.

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

Why does this matter?

Why does this matter?

Philosophical remark

Since Newton's method is affine covariant, *the conditions for any theorem guaranteeing its convergence* should also be affine covariant.
Why does this matter?

Philosophical remark

Since Newton's method is affine covariant, *the conditions for any theorem guaranteeing its convergence* should also be affine covariant.

This is not true of proofs found in many books!

Why does this matter?

Philosophical remark

Since Newton's method is affine covariant, *the conditions for any theorem guaranteeing its convergence* should also be affine covariant.

This is not true of proofs found in many books!

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

 Invented linear programming when consulting for the Leningrad Plywood Trust.



Leonid Kantorovich (1912–1986).

- Invented linear programming when consulting for the Leningrad Plywood Trust.
- Instrumental in saving millions of lives during the siege of Leningrad.



Leonid Kantorovich (1912–1986).

- Invented linear programming when consulting for the Leningrad Plywood Trust.
- Instrumental in saving millions of lives during the siege of Leningrad.
- Involved in the Soviet nuclear bomb project.



Leonid Kantorovich (1912-1986).

- Invented linear programming when consulting for the Leningrad Plywood Trust.
- Instrumental in saving millions of lives during the siege of Leningrad.
- Involved in the Soviet nuclear bomb project.
- Nearly sent to the gulag for "shadow prices".



Leonid Kantorovich (1912-1986).

- Invented linear programming when consulting for the Leningrad Plywood Trust.
- Instrumental in saving millions of lives during the siege of Leningrad.
- Involved in the Soviet nuclear bomb project.
- Nearly sent to the gulag for "shadow prices".
- Pseudo-Nobel prize in Economics (1975).



Leonid Kantorovich (1912–1986).



Kantorovich's theorem (1948) is an example of the application of the Banach contraction mapping theorem. It does not assume the existence of a solution: given certain conditions on the residual and initial guess, it proves the existence and local uniqueness of a solution, and the convergence of the Newton iteration.

Kantorovich's theorem (1948) is an example of the application of the Banach contraction mapping theorem. It does not assume the existence of a solution: given certain conditions on the residual and initial guess, it proves the existence and local uniqueness of a solution, and the convergence of the Newton iteration.

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) in finite dimensions)

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 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 (1) $\|DF(\mathbf{x}_0)^{-1}F(\mathbf{x}_0)\| \leq \frac{r}{2}$,

Theorem (Kantorovich (1948) in finite dimensions)

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 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

(1)
$$||DF(\mathbf{x}_0)^{-1}F(\mathbf{x}_0)|| \le \frac{r}{2}$$
,
(2) For all $\tilde{\mathbf{x}}, \mathbf{x} \in B(\mathbf{x}_0, r)$,

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

Then

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

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}^* of *F*.

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}^* of F. (3) For each $i \ge 0$,

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

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}^* of F. (3) For each $i \ge 0$,

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

(4) The root \mathbf{x}^* is locally unique, i.e. \mathbf{x}^* 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

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

Bonus: The Davidenko differential equation

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

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

Philosophical question

Is there a $\mathit{curve}\;\mathbf{x}(s),s\in[0,\infty),$ associated with this sequence?

Bonus: The Davidenko differential equation

Newton's method applied to $F(\mathbf{x}) = \mathbf{0}$ produces a sequence $\mathbf{x}_0, \mathbf{x}_1, \mathbf{x}_2, \dots, \quad \mathbf{x}_i \in \mathbb{R}^N.$

Philosophical question

Is there a curve $\mathbf{x}(s), s \in [0,\infty)$, associated with this sequence?

Yes. The Davidenko differential equation is

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



Victor Davidenko, 1914-1983

Bonus: The Davidenko differential equation

Newton's method applied to $F(\mathbf{x}) = \mathbf{0}$ produces a sequence $\mathbf{x}_0, \mathbf{x}_1, \mathbf{x}_2, \dots, \quad \mathbf{x}_i \in \mathbb{R}^N.$

Philosophical question

Is there a curve $\mathbf{x}(s), s \in [0,\infty)$, associated with this sequence?

Yes. The Davidenko differential equation is

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

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

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



Victor Davidenko, 1914–1983

Theorem

For any $\mathbf{x}_0 \in \mathbb{R}^N$, the solution curve of the Davidenko differential equation ends either at

 \blacktriangleright a root \mathbf{x}^{\star} , or at

Theorem

For any $\mathbf{x}_0 \in \mathbb{R}^N$, the solution curve of the Davidenko differential equation ends either at

- \blacktriangleright a root \mathbf{x}^{\star} , or at
- \blacktriangleright a singular point of DF.

Theorem

For any $\mathbf{x}_0 \in \mathbb{R}^N$, the solution curve of the Davidenko differential equation ends either at

- \blacktriangleright a root \mathbf{x}^{\star} , or at
- \blacktriangleright a singular point of DF.

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*.

Theorem

For any $\mathbf{x}_0 \in \mathbb{R}^N$, the solution curve of the Davidenko differential equation ends either at

- \blacktriangleright a root $\mathbf{x}^{\star},$ or at
- \blacktriangleright a singular point of DF.

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*.

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

Section 8

Newton fractals

Consider the problem

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

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

Consider the problem

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

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

We know this has three solutions,

$$z = 1$$
, $z = -1/2 + i\sqrt{3}/2$, 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,

Consider the problem

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

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

We know this has three solutions,

$$z = 1$$
, $z = -1/2 + i\sqrt{3}/2$, 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,

Consider the problem

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

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

We know this has three solutions,

$$z = 1$$
, $z = -1/2 + i\sqrt{3}/2$, 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 $z^3 - 1 = 0$.

P. E. Farrell (Oxford)



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

P. E. Farrell (Oxford)

Some useful websites:

- https://attr.actor/snapshots/dxhdzbzwmylmtywj
- https://newtonfractal.starfree.app/
- https://www.youtube.com/watch?v=-RdOwhmqP5s
Section 9

Algorithms for optimisation problems

Optimisation is fundamental to applied mathematics and engineering. It is also the engine that powers machine learning.

Optimisation is fundamental to applied mathematics and engineering. It is also the engine that powers machine learning.

Nature optimizes. Physical systems tend to a state of minimum energy. The molecules in an isolated chemical system react with each other until the total potential energy of their electrons is minimized. Rays of light follow paths that minimize their travel time.

Optimisation is fundamental to applied mathematics and engineering. It is also the engine that powers machine learning.

Nature optimizes. Physical systems tend to a state of minimum energy. The molecules in an isolated chemical system react with each other until the total potential energy of their electrons is minimized. Rays of light follow paths that minimize their travel time.

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.

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

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

We assume f is bounded below (e.g. $f(x) = -x^2$ has no min over \mathbb{R}).

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

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

We assume f is bounded below (e.g. $f(x) = -x^2$ has no min over \mathbb{R}).

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

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

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

We assume f is bounded below (e.g. $f(x) = -x^2$ has no min over \mathbb{R}).

We want an argument \mathbf{x}^* that satisfies $f(\mathbf{x}^*) \leq f(\mathbf{x})$ 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}^* such that there is a neighbourhood \mathcal{N} around \mathbf{x}^* so that

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

In our case, the optimality conditions are that the gradient $g : \mathbb{R}^N \to \mathbb{R}^N$ is zero at a local minimiser:

$$g(\mathbf{x}^{\star}) \coloneqq \nabla f(\mathbf{x}^{\star}) = Df(x^{\star})^{\top} = \begin{pmatrix} \frac{\partial f}{\partial \mathbf{x}^{1}}(\mathbf{x}^{\star})\\ \vdots\\ \frac{\partial f}{\partial \mathbf{x}^{N}}(\mathbf{x}^{\star}) \end{pmatrix} = \mathbf{0}.$$

0.0

In our case, the optimality conditions are that the gradient $g : \mathbb{R}^N \to \mathbb{R}^N$ is zero at a local minimiser:

$$g(\mathbf{x}^{\star}) \coloneqq \nabla f(\mathbf{x}^{\star}) = Df(x^{\star})^{\top} = \begin{pmatrix} \frac{\partial f}{\partial \mathbf{x}^{1}}(\mathbf{x}^{\star})\\ \vdots\\ \frac{\partial f}{\partial \mathbf{x}^{N}}(\mathbf{x}^{\star}) \end{pmatrix} = \mathbf{0}.$$

In other words, if we can find roots of g, we can find local minima of f!

In our case, the optimality conditions are that the gradient $g : \mathbb{R}^N \to \mathbb{R}^N$ is zero at a local minimiser:

$$g(\mathbf{x}^{\star}) \coloneqq \nabla f(\mathbf{x}^{\star}) = Df(x^{\star})^{\top} = \begin{pmatrix} \frac{\partial f}{\partial \mathbf{x}^{1}}(\mathbf{x}^{\star})\\ \vdots\\ \frac{\partial f}{\partial \mathbf{x}^{N}}(\mathbf{x}^{\star}) \end{pmatrix} = \mathbf{0}.$$

In other words, if we can find roots of g, we can find local minima of f!

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

global minimisers \subset local minimisers \subset critical points.

```
global minimisers \subset local minimisers \subset critical points.
```

We will have to be careful, when looking for critical points, to find only the local minimisers we're interested in.

```
global minimisers \subset local minimisers \subset critical points.
```

We will have to be careful, when looking for critical points, to find only the local minimisers we're interested in.

Local minimisers can be distinguished by studying the *second-order sufficiency conditions*. We won't see these.

```
global minimisers \subset local minimisers \subset critical points.
```

We will have to be careful, when looking for critical points, to find only the local minimisers we're interested in.

Local minimisers can be distinguished by studying the *second-order sufficiency conditions*. We won't see these.

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:

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

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

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

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



Albert Tucker, 1905-1995

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

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

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.

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



William Karush, 1917-1997



Harold Kuhn, 1925-2014



Albert Tucker, 1905-1995

Section 10

Newton's method for optimisation

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

.

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

The Hessian is always symmetric for $f \in C^2(\mathbb{R}^N, \mathbb{R})$.

.

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

The Hessian is always symmetric for $f \in C^2(\mathbb{R}^N, \mathbb{R})$.

Applying Newton's method to find roots of $\nabla f(x)$, we get

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

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

The Hessian is always symmetric for $f \in C^2(\mathbb{R}^N, \mathbb{R})$.

Applying Newton's method to find roots of $\nabla f(x)$, we get

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

Suppose we're at iterate x_i and we'd like to minimise f. We don't know how, so we'll replace f with a local quadratic model:

$$f(\mathbf{x}_i + \delta \mathbf{x}) \approx m(\delta \mathbf{x}) \coloneqq f(\mathbf{x}_i) + \nabla f(\mathbf{x}_i)^\top \delta \mathbf{x} + \frac{1}{2} \delta \mathbf{x}^\top H f(\mathbf{x}_i) \delta \mathbf{x}.$$

Suppose we're at iterate x_i and we'd like to minimise f. We don't know how, so we'll replace f with a local quadratic model:

$$f(\mathbf{x}_i + \delta \mathbf{x}) \approx m(\delta \mathbf{x}) \coloneqq f(\mathbf{x}_i) + \nabla f(\mathbf{x}_i)^\top \delta \mathbf{x} + \frac{1}{2} \delta \mathbf{x}^\top H f(\mathbf{x}_i) \delta \mathbf{x}.$$

We can decide what the update $\delta {\bf x}$ should be by solving $\nabla m(\delta {\bf x})=0,$ which yields the update

$$\delta \mathbf{x} = -[Hf(\mathbf{x}_i)]^{-1} \nabla f(\mathbf{x}_i).$$

Suppose we're at iterate x_i and we'd like to minimise f. We don't know how, so we'll replace f with a local quadratic model:

$$f(\mathbf{x}_i + \delta \mathbf{x}) \approx m(\delta \mathbf{x}) \coloneqq f(\mathbf{x}_i) + \nabla f(\mathbf{x}_i)^\top \delta \mathbf{x} + \frac{1}{2} \delta \mathbf{x}^\top H f(\mathbf{x}_i) \delta \mathbf{x}.$$

We can decide what the update $\delta {\bf x}$ should be by solving $\nabla m(\delta {\bf x})=0,$ which yields the update

$$\delta \mathbf{x} = -[Hf(\mathbf{x}_i)]^{-1} \nabla f(\mathbf{x}_i).$$

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

Section 11

Quasi-Newton methods

1. We only want to find local *minimisers*.

1. We only want to find local *minimisers*.

We often have huge $N \gg 1$, causing two more difficulties:

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

1. We only want to find local *minimisers*.

We often have huge $N \gg 1$, causing two more difficulties:

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

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 $\mathcal{O}(N^3)$.

1. We only want to find local *minimisers*.

We often have huge $N \gg 1$, causing two more difficulties:

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

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 $\mathcal{O}(N^3)$.

But for many problems no such nice structure exists (e.g. neural networks).

1. We only want to find local *minimisers*.

We often have huge $N \gg 1$, causing two more difficulties:

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

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 $\mathcal{O}(N^3)$.

But for many problems no such nice structure exists (e.g. neural networks).

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 :

1. $B_i = Hf(\mathbf{x}_i)$. Newton again. Quadratic convergence, often impractical.
Here are some choices for B_i :

- 1. $B_i = Hf(\mathbf{x}_i)$. Newton again. Quadratic convergence, often impractical.
- 2. $B_i = I \in \mathbb{R}^{N \times N}$. Gradient descent. Linear convergence, very slow.

Here are some choices for B_i :

- 1. $B_i = Hf(\mathbf{x}_i)$. Newton again. Quadratic convergence, often impractical.
- 2. $B_i = I \in \mathbb{R}^{N \times N}$. Gradient descent. Linear convergence, very slow.

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



Here are some choices for B_i :

- 1. $B_i = Hf(\mathbf{x}_i)$. Newton again. Quadratic convergence, often impractical.
- 2. $B_i = I \in \mathbb{R}^{N \times N}$. Gradient descent. Linear convergence, very slow.

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



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

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

which is the higher-order generalisation of the secant method.

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

which is the higher-order generalisation of the secant method.

In one dimension, this secant condition is enough to approximate $f''(x_i)$. But in higher dimensions it is not; we have N equations, but N(N+1)/2 variables to define B_{i+1} . So how do we fill in the missing information?

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

which is the higher-order generalisation of the secant method.

In one dimension, this secant condition is enough to approximate $f''(x_i)$. But in higher dimensions it is not; we have N equations, but N(N+1)/2 variables to define B_{i+1} . So how do we fill in the missing information?

BFGS proposed to choose, among all symmetric matrices satisfying the secant condition, the one whose inverse is closest to B_i^{-1} :

$$\begin{split} B_{i+1} &= \underset{B \in \mathbb{R}^{N \times N}}{\operatorname{argmin}} \qquad \|B^{-1} - B_i^{-1}\| \\ &\text{subject to} \quad B = B^\top, \\ &\quad B(\mathbf{x}_{i+1} - \mathbf{x}_i) = \nabla f(\mathbf{x}_{i+1}) - \nabla f(\mathbf{x}_i). \end{split}$$

This means we now need to supply B_0 .

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

which is the higher-order generalisation of the secant method.

In one dimension, this secant condition is enough to approximate $f''(x_i)$. But in higher dimensions it is not; we have N equations, but N(N+1)/2 variables to define B_{i+1} . So how do we fill in the missing information?

BFGS proposed to choose, among all symmetric matrices satisfying the secant condition, the one whose inverse is closest to B_i^{-1} :

$$\begin{split} B_{i+1} &= \underset{B \in \mathbb{R}^{N \times N}}{\operatorname{argmin}} \quad \|B^{-1} - B_i^{-1}\| \\ &\text{subject to} \quad B = B^\top, \\ &\quad B(\mathbf{x}_{i+1} - \mathbf{x}_i) = \nabla f(\mathbf{x}_{i+1}) - \nabla f(\mathbf{x}_i). \end{split}$$

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} .

 $f(\mathbf{x}_{i+1}) < f(\mathbf{x}_i).$

 $f(\mathbf{x}_{i+1}) < f(\mathbf{x}_i).$

This will be achieved by ensuring that B_i is *positive definite*.

 $f(\mathbf{x}_{i+1}) < f(\mathbf{x}_i).$

This will be achieved by ensuring that B_i is positive definite.

Definition (positive-definite)

A matrix $A \in \mathbb{R}^{N \times N}$ is said to be positive-definite if $\mathbf{x}^{\top} A \mathbf{x} > 0$ for all nonzero $\mathbf{x} \in \mathbb{R}^{N}$.

 $f(\mathbf{x}_{i+1}) < f(\mathbf{x}_i).$

This will be achieved by ensuring that B_i is *positive definite*.

Definition (positive-definite)

A matrix $A \in \mathbb{R}^{N \times N}$ is said to be positive-definite if $\mathbf{x}^{\top} A \mathbf{x} > 0$ for all nonzero $\mathbf{x} \in \mathbb{R}^{N}$. This is equivalent to all of its eigenvalues being positive.

 $f(\mathbf{x}_{i+1}) < f(\mathbf{x}_i).$

This will be achieved by ensuring that B_i is positive definite.

Definition (positive-definite)

A matrix $A \in \mathbb{R}^{N \times N}$ is said to be positive-definite if $\mathbf{x}^{\top} A \mathbf{x} > 0$ for all nonzero $\mathbf{x} \in \mathbb{R}^{N}$. This is equivalent to all of its eigenvalues being positive.

Diagonal matrices

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

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

 $f(\mathbf{x}_{i+1}) < f(\mathbf{x}_i).$

This will be achieved by ensuring that B_i is *positive definite*.

Definition (positive-definite)

A matrix $A \in \mathbb{R}^{N \times N}$ is said to be positive-definite if $\mathbf{x}^{\top} A \mathbf{x} > 0$ for all nonzero $\mathbf{x} \in \mathbb{R}^{N}$. This is equivalent to all of its eigenvalues being positive.

Diagonal matrices

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

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

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

P. E. Farrell (Oxford)

$$f(\mathbf{x}_{i+1}) < f(\mathbf{x}_i)$$

we modify the iteration

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

to use a line search.

$$f(\mathbf{x}_{i+1}) < f(\mathbf{x}_i)$$

we modify the iteration

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

to use a *line search*.

The basic idea is this. The direction $\mathbf{d}_i = -B_i^{-1} \nabla f(\mathbf{x}_i)$ might point towards a minimum, but we may overshoot if $\|\mathbf{d}_i\|$ gets too large. We fix this by adjusting the magnitude of the step.

$$f(\mathbf{x}_{i+1}) < f(\mathbf{x}_i)$$

we modify the iteration

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

to use a *line search*.

The basic idea is this. The direction $\mathbf{d}_i = -B_i^{-1} \nabla f(\mathbf{x}_i)$ might point towards a minimum, but we may overshoot if $\|\mathbf{d}_i\|$ gets too large. We fix this by adjusting the magnitude of the step.

Define

$$\phi_i(t) := f(\mathbf{x}_i + t\mathbf{d}_i)$$

$$f(\mathbf{x}_{i+1}) < f(\mathbf{x}_i)$$

we modify the iteration

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

to use a *line search*.

The basic idea is this. The direction $\mathbf{d}_i = -B_i^{-1} \nabla f(\mathbf{x}_i)$ might point towards a minimum, but we may overshoot if $\|\mathbf{d}_i\|$ gets too large. We fix this by adjusting the magnitude of the step.

Define

$$\phi_i(t) := f(\mathbf{x}_i + t\mathbf{d}_i)$$

$$\phi_i'(0) = \nabla f(\mathbf{x}_i + 0\mathbf{d}_i)^\top \mathbf{d}_i$$

$$f(\mathbf{x}_{i+1}) < f(\mathbf{x}_i)$$

we modify the iteration

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

to use a *line search*.

The basic idea is this. The direction $\mathbf{d}_i = -B_i^{-1} \nabla f(\mathbf{x}_i)$ might point towards a minimum, but we may overshoot if $\|\mathbf{d}_i\|$ gets too large. We fix this by adjusting the magnitude of the step.

Define

$$\phi_i(t) := f(\mathbf{x}_i + t\mathbf{d}_i)$$

$$\phi_i'(0) = \nabla f(\mathbf{x}_i + 0\mathbf{d}_i)^\top \mathbf{d}_i$$
$$= -\nabla f(\mathbf{x}_i)^T B_i^{-1} \nabla f(\mathbf{x}_i)$$

$$f(\mathbf{x}_{i+1}) < f(\mathbf{x}_i)$$

we modify the iteration

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

to use a *line search*.

The basic idea is this. The direction $\mathbf{d}_i = -B_i^{-1} \nabla f(\mathbf{x}_i)$ might point towards a minimum, but we may overshoot if $\|\mathbf{d}_i\|$ gets too large. We fix this by adjusting the magnitude of the step.

Define

$$\phi_i(t) := f(\mathbf{x}_i + t\mathbf{d}_i)$$

$$\phi_i'(0) = \nabla f(\mathbf{x}_i + 0\mathbf{d}_i)^\top \mathbf{d}_i$$

= $-\nabla f(\mathbf{x}_i)^T B_i^{-1} \nabla f(\mathbf{x}_i)$
< 0.

Since $\phi'_i(0) < 0$, this means that there exists t > 0 such that

$$\phi_i(t) = f(\mathbf{x}_i + t\mathbf{d}_i) < f(\mathbf{x}_i) = \phi_i(0)$$

Since $\phi'_i(0) < 0$, this means that there exists t > 0 such that

$$\phi_i(t) = f(\mathbf{x}_i + t\mathbf{d}_i) < f(\mathbf{x}_i) = \phi_i(0)$$

so if we take a small enough step we will decrease f.

Since $\phi'_i(0) < 0$, this means that there exists t > 0 such that

$$\phi_i(t) = f(\mathbf{x}_i + t\mathbf{d}_i) < f(\mathbf{x}_i) = \phi_i(0)$$

so if we take a small enough step we will decrease f.

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 $\phi(t)$.

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

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

This is the *Rosenbrock* function and has unique minimiser $(x, y)^* = (1, 1)$.

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

This is the *Rosenbrock* function and has unique minimiser $(x, y)^* = (1, 1)$.

We solve this from $(x_0, y_0) = (-1.2, 1)^{\top}$ with gradient descent, Newton's method, and BFGS, with a Wolfe line search, until $\|\nabla f(x)\| < 10^{-5}$.

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

This is the *Rosenbrock* function and has unique minimiser $(x, y)^* = (1, 1)$.

We solve this from $(x_0, y_0) = (-1.2, 1)^{\top}$ with gradient descent, Newton's method, and BFGS, with a Wolfe line search, until $\|\nabla f(x)\| < 10^{-5}$.

| Gradient descent | Newton's method | BFGS |
|------------------------|-----------------------|-----------------------|
| 1.827×10^{-4} | 3.48×10^{-2} | 1.70×10^{-3} |
| 1.826×10^{-4} | 1.44×10^{-2} | 1.17×10^{-3} |
| 1.824×10^{-4} | 1.82×10^{-4} | 1.34×10^{-4} |
| 1.823×10^{-4} | 1.17×10^{-8} | 1.01×10^{-6} |

 $\|(x,y)-(x,y)^\star\|$ for the last 4 iterations.

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

This is the Rosenbrock function and has unique minimiser $(x, y)^{\star} = (1, 1)$.

We solve this from $(x_0, y_0) = (-1.2, 1)^{\top}$ with gradient descent, Newton's method, and BFGS, with a Wolfe line search, until $\|\nabla f(x)\| < 10^{-5}$.

| Gradient descent | Newton's method | BFGS |
|------------------------|-----------------------|-----------------------|
| 1.827×10^{-4} | 3.48×10^{-2} | 1.70×10^{-3} |
| 1.826×10^{-4} | 1.44×10^{-2} | 1.17×10^{-3} |
| 1.824×10^{-4} | 1.82×10^{-4} | 1.34×10^{-4} |
| 1.823×10^{-4} | 1.17×10^{-8} | 1.01×10^{-6} |

 $\|(x,y)-(x,y)^\star\|$ for the last 4 iterations.

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