

Document: Math 226A (Fall 2011)

Professor: Freund

Latest Update: April 2, 2012

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1 9-23-11

1.1 Announcements

- <http://www.math.ucdavis.edu/~freund/226A>
- Office Hours are Monday 12:30-2:30 at MSB 2140
- 5 homeworks
- Grading scheme
 - HW 50% (10% each)
 - Final 50% (open book, open notes)
 - * 226B and 226C will have final projects
- No textbook, several reference books are listed on the course webpage

1.2 General Remarks

Definition 1.1. Numerical Analysis

Numerical analysis is the study of algorithms for the problems of continuous mathematics (Trefethen).

Problem:

Input → Output: Solution(s) or No Solution

Numerical Methods:

(Approximate) Input → (Approximate) Solution

Sources of Errors:

1. Approximate input errors (cannot represent the data exactly, e.g. irrational numbers) ⇒ *conditioning* of the problem. These problems are inherent to the problem.
2. Rounding errors ⇒ *stability* of the algorithm
3. Approximation errors (i.e. terminating your algorithm when the result is “close enough”) ⇒ *convergence*

Example 1.2.

1. Systems of linear equations

$$Ax = b, \quad A \in \mathbb{C}^{n \times n}, b \in \mathbb{C}^n, A \text{ nonsingular}$$

Input: $A, b \rightarrow$ Solution: $x = A^{-1}b$ (obtained via LU factorization = modern Gaussian elimination)

2. Systems of nonlinear equations

$$\mathbf{f}(\mathbf{x}) = \mathbf{0}, \quad \text{where } \mathbf{f}(\mathbf{x}) = \mathbf{A}\mathbf{x} - \mathbf{b}$$

Input: $\mathbf{f} : D \rightarrow \mathbb{R}^n, D \subset \mathbb{R}^n$.

Solutions: $\mathbf{x} \in D$ such that $\mathbf{f}(\mathbf{x}) = \mathbf{0}$.

Standard method for solving: Newton's method

1.3 Introduction to LU Factorization

Remark 1.3. LU Factorization

Given: (nonsingular) $A \in \mathbb{C}^{n \times n}, b \in \mathbb{C}^n$.

Goal: solve $Ax = b$.

Gaussian elimination:

$$Ax = b \Leftrightarrow Ux = c, \quad \text{where } U = \begin{bmatrix} * & * & \dots & * \\ 0 & * & \dots & * \\ \vdots & \ddots & \ddots & * \\ 0 & \dots & 0 & * \end{bmatrix}$$

i.e. U is *upper-triangular*. How do we do this?

$$L_{n-1} \dots L_2 L_1 A = U \Leftrightarrow A = \underbrace{L_1^{-1} \dots L_{n-2}^{-1} L_{n-1}^{-1}}_L U$$
$$L = \begin{bmatrix} 1 & 0 & \dots & 0 \\ * & 1 & \dots & 0 \\ \vdots & \ddots & \ddots & 0 \\ * & \dots & * & 1 \end{bmatrix}$$

i.e. L is *lower-triangular*.

Example 1.4. General LU Factorization Example

$n = 4$

$$A = \begin{bmatrix} * & * & * & * \\ * & * & * & * \\ * & * & * & * \\ * & * & * & * \end{bmatrix} \xrightarrow{L_1} \begin{bmatrix} * & * & * & * \\ 0 & \# & \# & \# \\ 0 & \# & \# & \# \\ 0 & \# & \# & \# \end{bmatrix} = L_1 A$$

$$\xrightarrow{L_2} \begin{bmatrix} * & * & * & * \\ 0 & * & * & * \\ 0 & 0 & \# & \# \\ 0 & 0 & \# & \# \end{bmatrix} = L_2 L_1 A$$

$$\xrightarrow{L_3} \begin{bmatrix} * & * & * & * \\ 0 & * & * & * \\ 0 & 0 & * & * \\ 0 & 0 & 0 & \# \end{bmatrix} = L_3 L_2 L_1 A = U$$

*s represent unchanged nonzero entries, #s represent changed nonzero entries.

The L_i are Frobenius matrices.

Definition 1.5. Frobenius matrix

http://en.wikipedia.org/wiki/Frobenius_matrix

A *Frobenius matrix* is a square matrix with the following properties:

- all entries on the main diagonal are ones
- the entries below the main diagonal of at most one column are arbitrary
- every other entry is zero

Frobenius matrices are invertible. The following is an example:

$$A = \begin{bmatrix} 1 & 0 & 0 & \cdots & 0 \\ 0 & 1 & 0 & \cdots & 0 \\ 0 & a_{32} & 1 & \cdots & 0 \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ 0 & a_{n2} & 0 & \cdots & 1 \end{bmatrix}, \quad A^{-1} = \begin{bmatrix} 1 & 0 & 0 & \cdots & 0 \\ 0 & 1 & 0 & \cdots & 0 \\ 0 & -a_{32} & 1 & \cdots & 0 \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ 0 & -a_{n2} & 0 & \cdots & 1 \end{bmatrix}$$

2 9-26-11

2.1 LU Factorization (Continued)

Example 2.1.

$$\begin{aligned}
 A &= \begin{bmatrix} 2 & 1 & 1 & 0 \\ 4 & 3 & 3 & 1 \\ 8 & 7 & 9 & 5 \\ 6 & 7 & 9 & 8 \end{bmatrix} \\
 L_1 A &= \begin{bmatrix} 1 & 0 & 0 & 0 \\ -2 & 1 & 0 & 0 \\ -4 & 0 & 1 & 0 \\ -3 & 0 & 0 & 1 \end{bmatrix} A = \begin{bmatrix} 2 & 1 & 1 & 0 \\ 0 & 1 & 1 & 1 \\ 0 & 3 & 5 & 5 \\ 0 & 4 & 6 & 8 \end{bmatrix} \\
 L_2 L_1 A &= \begin{bmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & -3 & 1 & 0 \\ 0 & -4 & 0 & 1 \end{bmatrix} \begin{bmatrix} 2 & 1 & 1 & 0 \\ 0 & 1 & 1 & 1 \\ 0 & 3 & 5 & 5 \\ 0 & 4 & 6 & 8 \end{bmatrix} = \begin{bmatrix} 2 & 1 & 1 & 0 \\ 0 & 1 & 1 & 1 \\ 0 & 0 & 2 & 2 \\ 0 & 0 & 2 & 4 \end{bmatrix} \\
 L_3 L_2 L_1 &= \begin{bmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & -1 & 1 \end{bmatrix} L_2 L_1 A = \begin{bmatrix} 2 & 1 & 1 & 0 \\ 0 & 1 & 1 & 1 \\ 0 & 0 & 2 & 2 \\ 0 & 0 & 0 & 2 \end{bmatrix} = U \quad \text{upper-triangular} \\
 A &= LU \\
 L &= L_1^{-1} L_2^{-1} L_3^{-1} = \begin{bmatrix} 1 & 0 & 0 & 0 \\ 2 & 1 & 0 & 0 \\ 4 & 0 & 1 & 0 \\ 3 & 0 & 0 & 1 \end{bmatrix} \begin{bmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 3 & 1 & 0 \\ 0 & 4 & 0 & 1 \end{bmatrix} \begin{bmatrix} 1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 1 \end{bmatrix} \\
 &= \begin{bmatrix} 1 & 0 & 0 & 0 \\ 2 & 1 & 0 & 0 \\ 4 & 3 & 1 & 0 \\ 3 & 4 & 1 & 1 \end{bmatrix} \quad \text{unit lower-triangular}
 \end{aligned}$$

Remark 2.2. Convention (Matlab Notation)

$$\begin{aligned}
 \mathbf{A} &= [a_{jk}]_{j,k=1,2,\dots,n} \in \mathbb{C}^{n \times n} \\
 a_{j,k:n} &= [a_{j,k} \ a_{j,k+1} \ \cdots \ a_{j,n}]
 \end{aligned}$$

Remark 2.3. Algorithm: LU Factorization without pivoting

Input: $\mathbf{A} \in \mathbb{C}^{n \times n}$

Output: L, U such that $A = LU$

Set $U = A$, $L = I$ ($n \times n$ identity matrix).

- for $k = 1, 2, \dots, n - 1$
 - for $j = k + 1, k + 2, \dots, n$
 - * $l_{j,k} = u_{j,k}/u_{k,k}$ (potential problem)
 - * $u_{j,k:n} = u_{j,k:n} - l_{j,k}u_{k,k:n}$
 - end
- end

Output: $L, U \in \mathbb{C}^{n \times n}$ such that $A = LU$

Remark 2.4. Operation Count

Additions: $\sum_{k=1}^{n-1} (n-k)(n-k+1) = \sum_{l=1}^{n-1} l(l+1) \approx \frac{n^3}{3}$ (where $l = n - k$)

Multiplications: $\approx \frac{n^3}{3}$

Divisions: $\approx n^2 \Rightarrow$ ignore (because it is a lower order term)

Total Work: $\approx \frac{2n^3}{3}$ flops

Definition 2.5. Flop

A *flop* is a floating-point operation: addition, subtraction, multiplication, division, square root.

Remark 2.6. Solution of $Ax = b$

$$\underbrace{\mathbf{A}}_{=LU} \mathbf{x} = \mathbf{b} \quad (2.1)$$

$$L \underbrace{Ux}_{=y} = \mathbf{b} \quad (2.2)$$

$$Ly = \mathbf{b} \quad (2.1)$$

$$Ux = y \quad (2.2)$$

$Ly = b$ is easily solved via forward substitution. $Ux = y$ is easily solved by backward substitution.

One triangular solve requires $\approx n^2$ flops.

2.2 Problems with LU Factorization without Pivoting

LU factorization without pivoting is unstable!

Example 2.7. Error from LU Factorization without Pivoting

$$A = \begin{bmatrix} 10^{-20} & 1 \\ 1 & 1 \end{bmatrix}, \quad b = \begin{bmatrix} 1 \\ 0 \end{bmatrix}$$

With exact arithmetic:

$$L = \begin{bmatrix} 1 & 0 \\ 10^{20} & 1 \end{bmatrix}, \quad U = \begin{bmatrix} 10^{-20} & 1 \\ 0 & 1 - 10^{20} \end{bmatrix}$$

Now solve a system:

$$\begin{aligned} Ly &= b \\ y &= \begin{bmatrix} 1 \\ -10^{20} \end{bmatrix} \\ Ux &= y \\ x &= \frac{1}{1 - 10^{-20}} \begin{bmatrix} -1 \\ 1 \end{bmatrix} \\ &\approx \begin{bmatrix} -1 \\ 1 \end{bmatrix} \end{aligned}$$

In a computer:

Floating point arithmetic does everything accurately up to ≈ 16 digits.

$$\begin{aligned} \tilde{L} &= L = \begin{bmatrix} 1 & 0 \\ 10^{20} & 1 \end{bmatrix}, \quad \tilde{U} = \begin{bmatrix} 10^{-20} & 1 \\ 0 & -10^{20} \end{bmatrix} \\ \tilde{L}\tilde{U} &= \begin{bmatrix} 10^{-20} & 1 \\ 1 & 0 \end{bmatrix} \not\approx A \end{aligned}$$

If you use this to compute the answer, you will get a bogus answer.

3 9-28-11

3.1 Continued from 9-26-11...

Example 3.1. *Continuing from last time...*

$$\begin{aligned}\tilde{L}\tilde{y} &= b \\ \tilde{y} &= \begin{bmatrix} 1 \\ 10^{-20} \end{bmatrix} \\ \tilde{U}\tilde{x} &= \tilde{y} \\ \tilde{x} &= \begin{bmatrix} 0 \\ 1 \end{bmatrix} \neq x\end{aligned}$$

The culprit is 10^{20} .

3.2 Pivoting

$$\begin{aligned}A &= \begin{bmatrix} 10^{-20} & 1 \\ 1 & 1 \end{bmatrix} \\ P &= \begin{bmatrix} 0 & 1 \\ 1 & 0 \end{bmatrix} \quad (\text{permutation matrix}) \\ PA &= \begin{bmatrix} 1 & 1 \\ 20^{-20} & 1 \end{bmatrix} \\ &= \underbrace{\begin{bmatrix} 1 & 0 \\ 10^{-20} & 1 \end{bmatrix}}_{=L} \underbrace{\begin{bmatrix} 1 & 1 \\ 0 & 1 - 10^{-20} \end{bmatrix}}_{=U}\end{aligned}$$

Floating point arithmetic:

$$\begin{aligned}\tilde{L} &= L \\ \tilde{U} &= \begin{bmatrix} 1 & 1 \\ 0 & 1 \end{bmatrix} \\ PAx &= Pb \\ LUx &= \begin{bmatrix} 0 \\ 1 \end{bmatrix} \\ \tilde{L}\tilde{y} &= Pb \\ \tilde{y} &= \begin{bmatrix} 0 \\ 1 \end{bmatrix} \\ \tilde{U}\tilde{x} &= \begin{bmatrix} 0 \\ 1 \end{bmatrix} \\ \tilde{x} &= \begin{bmatrix} -1 \\ 1 \end{bmatrix} \approx x = A^{-1}b\end{aligned}$$

Every time you use Gaussian elimination you need to use pivoting.

Example 3.2. General Pivoting Example

$n = 4, k = 2$ (4×4 matrix)

$$U = \begin{bmatrix} * & * & * & * \\ 0 & (\ast) & * & * \\ 0 & * & * & * \\ 0 & * & * & * \end{bmatrix}$$

$$\xrightarrow{P_2} \begin{bmatrix} * & * & * & * \\ 0 & (\#) & \# & \# \\ 0 & * & * & * \\ 0 & \# & \# & \# \end{bmatrix} \quad \text{exchange rows 2 \& 4}$$

$$\xrightarrow{L_2} \begin{bmatrix} * & * & * & * \\ 0 & * & * & * \\ 0 & 0 & \# & \# \\ 0 & 0 & \# & \# \end{bmatrix}$$

In general:

$$L_{n-1}P_{n-1} \cdots L_3L_2P_2L_1P_1A = U$$

If you don't need to pivot, then use $P_i = I$ (the identity).

Example 3.3. Specific Pivoting Example

$$\begin{aligned}
 A &= \begin{bmatrix} 2 & 1 & 1 & 0 \\ 4 & 3 & 3 & 1 \\ 8 & 7 & 9 & 5 \\ 6 & 7 & 9 & 8 \end{bmatrix} \\
 \underbrace{\begin{bmatrix} 0 & 0 & 1 & 0 \\ 0 & 1 & 0 & 0 \\ 1 & 0 & 0 & 0 \\ 0 & 0 & 0 & 1 \end{bmatrix}}_{P_1} A &= \begin{bmatrix} 8 & 7 & 9 & 5 \\ 4 & 3 & 3 & 1 \\ 2 & 1 & 1 & 0 \\ 6 & 7 & 9 & 8 \end{bmatrix} \\
 L_1 &= \begin{bmatrix} 1 & 0 & 0 & 0 \\ -\frac{1}{2} & 1 & 0 & 0 \\ -\frac{1}{4} & 0 & 1 & 0 \\ -\frac{3}{4} & 0 & 0 & 1 \end{bmatrix} \\
 L_1 P_1 A &= \begin{bmatrix} 8 & 7 & 9 & 5 \\ 0 & \left(\frac{7}{4}\right) & \frac{9}{4} & \frac{17}{4} \\ 0 & -\frac{3}{4} & -\frac{5}{4} & -\frac{5}{4} \\ 0 & -\frac{1}{2} & -\frac{3}{2} & -\frac{3}{2} \end{bmatrix} \\
 &\vdots \\
 U &= \begin{bmatrix} 8 & 7 & 9 & 5 \\ 0 & \frac{7}{4} & \frac{9}{4} & \frac{17}{4} \\ 0 & 0 & -\frac{6}{7} & -\frac{2}{7} \\ 0 & 0 & 0 & \frac{2}{3} \end{bmatrix}
 \end{aligned}$$

Example 3.4. Pivoting in Matrix Form

$n = 4 :$

$$\begin{aligned}
 L_3 P_3 &\underbrace{I}_{=P_3^T P_3} L_2 P_2 L_1 & \underbrace{I}_{=P_2^T P_3^T P_3 P_2} P_1 A = U \\
 &= L'_3 \underbrace{(P_3 L_2 P_3^T)}_{=L'_2} \underbrace{(P_3 P_2 L_1 P_2^T P_3^T)}_{=L'_1} \underbrace{(P_3 P_2 P_1)}_{=P} A = U \\
 L'_3 L'_2 L'_1 (PA) &= U
 \end{aligned}$$

L'_1 ($\neq L_1^T$) has the same structure as L_1 , with nontrivial elements in column 1 reordered.

$$L_1 = \begin{bmatrix} 1 & 0 & 0 & 0 \\ -\frac{1}{2} & 1 & 0 & 0 \\ -\frac{1}{4} & 0 & 1 & 0 \\ -\frac{3}{4} & 0 & 0 & 1 \end{bmatrix}, \quad L'_1 = \begin{bmatrix} 1 & 0 & 0 & 0 \\ * & 1 & 0 & 0 \\ * & 0 & 1 & 0 \\ * & 0 & 0 & 1 \end{bmatrix}$$

$$PA = \underbrace{(L'_1)^{-1}(L'_2)^{-1}(L'_3)^{-1}}_{=L} U$$

General case:

$$PA = LU$$

where P is the permutation matrix and L and U are lower and upper triangular matrices, respectively.

Remark 3.5. Algorithm: LU factorization with partial pivoting

Input: $A \in \mathbb{C}^{n \times n}$

Set $U = A$, $L = I$, $P = I$.

- for $k = 1, 2, \dots, n - 1$
 - choose $i \geq k$ such that $|u_{ik}| = \max_{k \leq l \leq n} |u_{lk}|$
 - if $u_{kk} = 0$, stop: A is singular
 - $u_{k,k:n} \leftrightarrow u_{i,k:n}$ (interchange rows i and k)
 - $l_{k,1:k-1} \leftrightarrow l_{i,1:k-1}$
 - $p_{k,:} \leftrightarrow p_{i,:}$
 - for $j = k + 1, k + 2, \dots, n$
 - * $l_{jk} = u_{jk}/u_{kk}$
 - * $u_{j,k:n} = u_{j,k:n} - l_{jk}u_{k,k:n}$
 - end (j)
- end (k)

Output: U, L, P such that

$$PA = LU$$

where $L = [l_{jk}]$ with $|l_{jk}| \leq 1$ for all $j > k$.

4 9-30-11

4.1 LU Factorization Recap

$$PA = LU$$

$$L = \begin{bmatrix} 1 & 0 & \cdots & 0 \\ l_{21} & 1 & \ddots & \vdots \\ \vdots & \ddots & \ddots & 0 \\ l_{n1} & \cdots & l_{n,n-1} & 1 \end{bmatrix}, \quad |l_{jk}| \leq 1$$

LU factorization with partial pivoting is stable in practice.

4.2 Newton's Method

Newton's method is not a finite method, meaning that you would have to run it infinitely many times in order for it to arrive at the answer.

Remark 4.1. Newton's Method Overview

Given: $f : D \rightarrow \mathbb{R}^n$, $D \subset \mathbb{R}^n$

$$f = \begin{bmatrix} f_1 \\ f_2 \\ \vdots \\ f_n \end{bmatrix}, \quad f_j : D \rightarrow \mathbb{R}, \quad j = 1, 2, \dots, n$$

$$x = \begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_n \end{bmatrix} \in D \subset \mathbb{R}^n, \quad x_j \in \mathbb{R}, \quad j = 1, 2, \dots, n$$

Goal: Find $x \in D$ such that

$$f(\mathbf{x}) = \mathbf{0}.$$

Thus, we have n nonlinear equations for n unknowns.

Example 4.2. Newton's Method Example 1: Possible Scenarios

$$D = \mathbb{R}^n, \quad f : \mathbb{R}^n \rightarrow \mathbb{R}^n,$$

$$\begin{aligned} f(\mathbf{x}) &= \mathbf{b} - \mathbf{Ax}, & \mathbf{b} \in \mathbb{R}^n, \quad \mathbf{A} \in \mathbb{R}^{n \times n} \\ f(\mathbf{x}) &= 0 \Leftrightarrow \mathbf{Ax} = \mathbf{b} \end{aligned}$$

There are 3 possibilities:

1. 1 solution
2. no solution
3. infinitely many solutions

Example 4.3. Newton's Method Example 2

$f = f_\alpha : \mathbb{R} \rightarrow \mathbb{R}$, $f_\alpha(x) = x^2 + \alpha$, $\alpha \in \mathbb{R}$ is a parameter. This is an upward facing parabola with minimum value α .

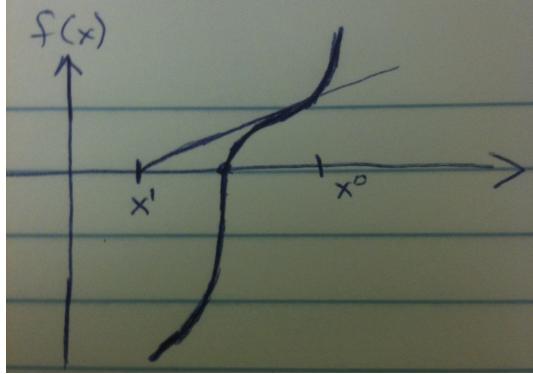
1. $f(x)$ has no solution if $\alpha > 0$
2. $f(x)$ has one solution if $\alpha = 0$ ($x = 0$)
3. $f(x)$ has two solutions if $\alpha < 0$ ($x = \pm\sqrt{-\alpha}$)

Example 4.4. Newton's Method Example 3

$n = 1$, $D = \mathbb{R}^+ = \{x \in \mathbb{R} \mid x > 0\}$, $f(x) = \sin \frac{1}{x}$. Thus, f oscillates faster as $x \rightarrow 0$.

$$\begin{aligned}f(x) = 0 &\Leftrightarrow \frac{1}{x} = j\pi, \quad j = 1, 2, \dots \\x &= \frac{1}{j\pi}\end{aligned}$$

Remark 4.5. Newton's Method Algorithm for $n = 1$



x^0 = initial guess

$$f(x^*) = 0$$

1. Linearize $f(x)$ around x^0 :

$$f(x) = f(x^0) + f'(x^0)(x - x^0) + O((x - x^0)^2) \equiv t(x)$$

2. Get new approximation x^1 by setting $t(x^1) = 0$ and solving for x^1 :

$$x^1 = x^0 - \frac{f(x^0)}{f'(x^0)}$$

(provided $f'(x) \neq 0$).

3. Repeat with x^0 replaced by x^1 . Iterative process:

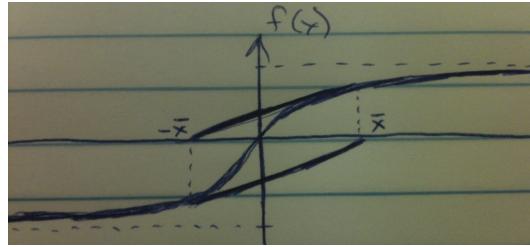
$$x^{k+1} = x^k - \frac{f(x^k)}{f'(x^k)}, \quad k = 0, 1, 2, \dots$$

Remark 4.6. Questions about Newton's Method

1. Can we guarantee that $\lim_{k \rightarrow \infty} x^k = x^*$?
2. If yes, what is the speed of convergence?

Example 4.7. Failure of Newton's Method

$$f : \mathbb{R} \rightarrow \mathbb{R}, f(x) = \tan^{-1} x, f'(x) = \frac{1}{1+x^2}.$$



f has a single zero: $x^* = 0$. We apply Newton's method. Choose $x^0 \in \mathbb{R}$.

$$x^{k+1} = x^k - \left(1 + (x^k)^2\right) \tan^{-1}(x^k), \quad k = 0, 1, 2, \dots$$

Convergence?

Let \bar{x} be the unique solution of

$$\begin{aligned} -\bar{x} &= \bar{x} - (1 + \bar{x}^2) \tan^{-1} \bar{x}, \quad \bar{x} > 0 \\ \bar{x} &= 1.39174 \dots \end{aligned}$$

where \bar{x} is computed by using Newton's Method for

$$g(x) = 2x - (1 + x^2) \tan^{-1} x.$$

Thus, if $x^0 = \bar{x}$, then $x^k = (-1)^k \bar{x}$, $k = 0, 1, \dots$

Similarly, if $x^0 = -\bar{x}$, then $x^{k+1} = (-1)^{k+1} \bar{x}$, $k = 0, 1, \dots$

If $|x^0| < \bar{x}$, then $\lim_{k \rightarrow \infty} x^k = 0 = x^*$

If $|x^0| > \bar{x}$, then $\lim_{k \rightarrow \infty} x^k = \infty$

5 10-3-11

5.1 Newton's Method (Continued)

Remark 5.1. *Newton's Method for $n \geq 1$*

$$f : D \rightarrow \mathbb{R}^n, \quad f = \begin{bmatrix} f_1 \\ f_2 \\ \vdots \\ f_n \end{bmatrix}, \quad x = \begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_n \end{bmatrix} \in D$$

Initial guess: $x^0 \in D$.

$$f(x) = f(x^0) + Df(x^0)(x - x^0) + O(\|x - x^0\|^2) \quad (5.1)$$

where

$$\|x\| = \sqrt{x_1^2 + x_2^2 + \cdots + x_n^2} = \text{Euclidean norm.}$$

$$Df(x) = \begin{bmatrix} \frac{\partial f_1}{\partial x_1} & \frac{\partial f_1}{\partial x_2} & \cdots & \frac{\partial f_1}{\partial x_n} \\ \vdots & \vdots & \ddots & \vdots \\ \frac{\partial f_n}{\partial x_1} & \frac{\partial f_n}{\partial x_2} & \cdots & \frac{\partial f_n}{\partial x_n} \end{bmatrix} = \left[\frac{\partial f_j}{\partial x_k} \right]_{j,k=1,2,\dots,n} = \text{Jacobian of } f$$

Obtain a new approximation x^1 by setting (5.1) to 0:

$$\begin{aligned} \mathbf{0} &= f(x^0) + Df(x^0) \underbrace{(x^1 - x^0)}_{\Delta x^0} \\ Df(x^0)\Delta x^0 &= -f(x^0) \\ x^1 &= x^0 + \Delta x^0. \end{aligned}$$

This is a system of n linear equations for n unknowns.

Remark 5.2. *Newton's Method Algorithm*

Choose $x^0 \in D$.

- for $k = 1, 2, \dots$ (open-ended), do:
 - if $Df(x^k)$ is singular, stop.
 - else, solve $Df(x^k)\Delta x^k = -f(x^k)$ for Δx^k
 - set $x^{k+1} = x^k + \Delta x^k$
 - if $x^{k+1} \notin D$, stop.
 - check for convergence: if $f(x^{k+1}) \approx 0$, stop.
- end (k)

Remark 5.3.

1. Need to solve a linear system at each k step \Rightarrow LU factorization of $A = Df(x^k)$

2. Newton's method is *affine-invariant*:

Let $M \in \mathbb{R}^{n \times n}$ be a fixed nonsingular matrix, and let $g : D \rightarrow \mathbb{R}^n$, $g(x) := Mf(x)$ for all $x \in D$. Thus, $f(x) = 0 \Leftrightarrow Mf(x) = \mathbf{0} \Leftrightarrow g(x) = 0$. Then Newton's method applied to f and g results in the same iterates x^k , $k = 0, 1, 2, \dots$

Proof. (That Newton's method is affine-invariant)

$$\begin{aligned} g(x) &= Mf(x) \\ Dg(x^k) &= MDf(x^k) \\ Dg(x^k)\Delta x^k &= -g(x^k) \quad \Leftrightarrow \quad Df(x^k)\Delta x^k = -f(x^k) \end{aligned}$$

□

5.2 Convergence of Newton's Method

Proposition 5.4. *Convergence*

Let x^* denote a zero of f , i.e. $f(x^*) = \mathbf{0}$. If $Df(x^*)$ is nonsingular and x^0 is “close” to x^* , then

$$\lim_{k \rightarrow \infty} x^k = x^*$$

and the speed of convergence is quadratic, i.e.

$$\|x^{k+1} - x^k\| \leq C \|x^k - x^*\|^2$$

Definition 5.5. *Quadratic Convergence*

The number of correct digits roughly doubles in each iteration.

Theorem 5.6. *Convergence of Newton's Method, Affine-Invariant Version*

Assumptions:

- $D \subset \mathbb{R}$ is convex
- $f : D \rightarrow \mathbb{R}^n$ is C^1
- $x^0 \in D$ and $Df(x^0)$ is nonsingular
- Df satisfies the *affine-invariant Lipschitz condition*:

$$\|(Df(x^0))^{-1}Df(y) - Df(x)\| \leq \gamma \|y - x\|$$

Here $\|\cdot\| : \mathbb{R}^n \rightarrow \mathbb{R}$ is a vector norm and $\|M\| := \max_{\|x\|=1} \|Mx\|$ is the associated matrix norm, which satisfies:

1. $\|Mx\| \leq \|M\|\|x\|$
2. $\|I\| = 1$
3. $\|MN\| \leq \|M\|\|N\|$

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6.1 Newton's Method Convergence Theorem

Theorem 6.1.

Given:

- $f : D \rightarrow \mathbb{R}^n$
- $x^0 \in D$ and $Df(x^0)$ is nonsingular
- $\|Df(x^0)^{-1}(Df(y) - Df(x))\| \leq \gamma \|y - x\|$ for all $x, y \in D$ (bigger $\gamma \Rightarrow$ more nonlinear)
- $\|Df(x^0)^{-1}f(x^0)\| \leq \alpha$ for some $\alpha > 0$
- $h := \alpha\gamma < \frac{1}{2}$

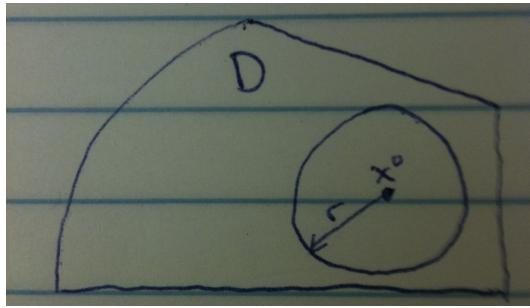


Figure 1: $S_r(x^0) := \{x \in \mathbb{R}^n \mid \|x - x^0\| < r\}$, $\overline{S_r(x^0)} := \{x \in \mathbb{R}^n \mid \|x - x^0\| \leq r\}$

- $\overline{S_r(x^0)} \subset D$, $r = \frac{1-\sqrt{1-2h}}{\gamma} > 0$

Claims:

1. All Newton iterates x^k satisfy $x^k \in S_r(x^0)$, $Df(x^k)$ is nonsingular
2. The iterates x^k converge quadratically to a zero of f , $x^* := \lim_{k \rightarrow \infty} x^k$, $x^* \in \overline{S_r(x^0)}$
3. x^* is the only zero of f in $D \cap S_R(x^0)$, $R := \frac{1+\sqrt{1-2h}}{\gamma}$

Note:

$$\|x^1 - x^0\| = \| - Df(x^0)^{-1}f(x^0) \| \leq \alpha = \frac{h}{\gamma}$$

Proof. (of parts 1 and 2 of the theorem)

We define two scalar sequences (h_k) and (γ_k) as

$$h_k := \begin{cases} h & k = 0 \\ \frac{h_{k-1}^2}{2(1-h_{k-1})^2} & k > 0 \end{cases}$$

$$\gamma_k := \begin{cases} \gamma & k = 0 \\ \frac{\gamma_{k-1}}{1-h_{k-1}} & k > 0 \end{cases}$$

It is easy to verify that for all $k = 0, 1, \dots$ we have

$$\begin{aligned}\lim_{j \rightarrow \infty} h_j &= 0 < h_{k+1} < h_k < \frac{1}{2} \\ \gamma_k < \gamma_{k+1} &< \frac{\gamma}{\sqrt{1-2h}} = \lim_{j \rightarrow \infty} \gamma_j\end{aligned}$$

2 Properties:

$$\begin{aligned}\frac{h_k}{\gamma_k} &\leq \frac{1}{2^k} \alpha \\ \frac{h_0}{\gamma_0} + \frac{h_1}{\gamma_1} + \dots + \frac{h_{k-1}}{\gamma_{k-1}} &= \frac{1}{\gamma} - \frac{1}{\gamma_k} < r = \frac{1 - \sqrt{1-2h}}{\gamma}\end{aligned}$$

Claim: For $k = 0, 1, \dots$, $x^k \in S_r(x^0)$, $Df(x^k)$ is nonsingular, and

$$\|Df(x^k)^{-1}(Df(y) - Df(x))\| \leq \gamma_k \|y - x\|, \quad x, y \in D.$$

We will also prove that

$$\|x^{k+1} - x^k\| \leq \frac{h_k}{\gamma_k} \leq \frac{1}{2^k} \alpha.$$

Proof by induction on k :

$k = 0 \Rightarrow$ satisfied (by our assumptions).

$k \geq 1$:

$$\begin{aligned}\|x^k - x^0\| &\leq \underbrace{\|x^k - x^{k-1}\|}_{\leq \frac{h_{k-1}}{\gamma_{k-1}}} + \underbrace{\|x^{k-1} - x^{k-2}\|}_{\leq \frac{h_{k-2}}{\gamma_{k-2}}} + \dots + \underbrace{\|x^1 - x^0\|}_{\leq \frac{h_0}{\gamma_0}} \\ &\leq \frac{h_0}{\gamma_0} + \frac{h_1}{\gamma_1} + \dots + \frac{h_{k-1}}{\gamma_{k-1}} \\ &< r\end{aligned}$$

Thus, $x^k \in S_r(x^0)$.

$$\begin{aligned}Df(x^k) &= \underbrace{Df(x^{k-1})}_{\text{nonsingular}} \left(\underbrace{I - Df(x^{k-1})^{-1}(Df(x^{k-1}) - Df(x^k))}_{:= A} \right) \tag{6.1} \\ \|A\| &= \|Df(x^{k-1})^{-1}(Df(x^k) - Df(x^{k-1}))\| \leq \gamma_{k-1} \underbrace{\|x^k - x^{k-1}\|}_{\leq \frac{h_{k-1}}{\gamma_{k-1}}} \\ &\leq h_{k-1} < \frac{1}{2} < 1\end{aligned}$$

Banach Lemma: If $\|A\| < 1$, then $I - A$ is nonsingular and

$$\|(I - A)^{-1}\| \leq \frac{1}{1 - h_{k-1}}.$$

Thus, by (6.1), $Df(x^k)$ is nonsingular and

$$\begin{aligned}\|Df(x^k)^{-1} \left[\underbrace{Df(x^{k-1})Df(x^{k-1})^{-1}}_I \right] Df(y) - Df(x)\| &\leq \underbrace{\|Df(x^k)^{-1}Df(x^{k-1})\|}_{(I-A)^{-1}} \|Df(x^{k-1})^{-1}(Df(y) - Df(x))\| \\ &\leq \frac{1}{1 - h_{k-1}} \gamma_{k-1} \|y - x\| \\ &= \gamma_k \|y - x\|, \quad x, y \in D.\end{aligned}$$

$$\begin{aligned}
x^{k+1} - x^k &= -Df(x^k)^{-1} \left(f(x^k) - \underbrace{f(x^{k-1}) - Df(x^{k-1})(x^k - x^{k-1})}_{=0} \right) \\
&= -Df(x^k)^{-1} \int_0^1 \left(Df(x^{k-1} + t(x^k - x^{k-1})) - Df(x^{k-1}) \right) (x^k - x^{k-1}) dt \\
\|x^{k+1} - x^k\| &\leq \int_0^1 \underbrace{\left\| Df(x^k)^{-1}(Df(x^{k-1} + t(x^k - x^{k-1})) - Df(x^{k-1})) \right\|} \left\| x^k - x^{k-1} \right\| dt \\
&\leq \gamma_k \int_0^1 \left\| x^{k-1} + t(x^k - x^{k-1}) - x^{k-1} \right\| \left\| x^k - x^{k-1} \right\| dt \\
&= \gamma_k \|x^k - x^{k-1}\|^2 \int_0^1 t dt \\
&= \frac{\gamma_k}{2} \underbrace{\|x^k - x^{k-1}\|^2}_{\leq \left(\frac{h_{k-1}}{\gamma_{k-1}}\right)^2} \leq \frac{\gamma_h}{2} \frac{h_{k-1}^2}{\gamma_{k-1}^2} \\
&= \frac{h_k}{\gamma_k}
\end{aligned}$$

To be continued...

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7.1 Continued from 10-5-11...

$$\begin{aligned}
x^k &\in S_r(x^0) \\
\|Df(x^k)^{-1}(Df(y) - Df(x))\| &\leq \gamma_k \|y - x\|, \quad x, y \in D \\
\|x^{k+1} - x^k\| &\leq \frac{h_k}{\gamma_k} \leq \frac{\alpha}{2^k}
\end{aligned}$$

$f : D \rightarrow \mathbb{R}^n$, $f(x^*) = 0$.

Claim: (x^k) is a Cauchy sequence.

Let $m > k$.

$$\begin{aligned}
\|x^m - x^k\| &\leq \|x^m - x^{m-1}\| + \|x^{m-1} - x^{m-2}\| + \cdots + \|x^{k+1} - x^k\| \\
&\leq \frac{\alpha}{2^{m-1}} + \frac{\alpha}{2^{m-2}} + \cdots + \frac{\alpha}{2^k} \\
&\leq \frac{\alpha}{2^k} \left(1 + \frac{1}{2} + \frac{1}{4} + \cdots + \frac{1}{2^{m-k-1}} \right) \\
&\leq \frac{\alpha}{2^k} \left(\underbrace{1 + \frac{1}{2} + \frac{1}{4} + \cdots}_{=2} \right) \\
&\leq \frac{\alpha}{2^{k-1}}
\end{aligned}$$

Thus, (x^k) is Cauchy and therefore convergent:

$$\lim_{k \rightarrow \infty} x^k = x^*, \quad x^* \in \overline{S_r(x^0)}.$$

Claim: x^* is a zero of f .

$$\begin{aligned}
f(x^k) &= -Df(x^k)(x^{k+1} - x^k) \\
\|f(x^k)\| &\leq \underbrace{\|Df(x^k)\|}_{\leq \max_{x \in \overline{S_r(x^0)}} \|Df(x)\| := M} \underbrace{\|x^{k+1} - x^k\|}_{\leq \frac{\alpha}{2^k}} \\
&\leq \frac{\alpha M}{2^k}
\end{aligned} \tag{7.1}$$

Thus,

$$0 = \lim_{k \rightarrow \infty} f(x^k) = f(\lim_{k \rightarrow \infty} x^k) = f(x^*)$$

Claim: the convergence is quadratic.

$$\begin{aligned}
x^{k+1} &= x^k - Df(x^k)^{-1} f(x^k) && \text{by (7.1)} \\
x^{k+1} - x^* &= x^k - x^* - Df(x^k)^{-1}(f(x^k) - \underbrace{f(x^*)}_{=0}) \\
&= \int_0^1 Df(x^k)^{-1}(Df(x^k) - Df(x^* + t(x^k - x^*))(x^k - x^*)) dt \\
\|x^{k+1} - x^*\| &\leq \|x^k - x^*\| \cdot \int_0^1 \|Df(x^k)^{-1}(Df(x^k) - Df(x^* + t(x^k - x^*)))\| dt \\
&\leq \|x^k - x^*\| \gamma_k \int_0^1 \underbrace{\|x^k - x^* - t(x^k - x^*)\|}_{=(1-t)\|x^k - x^*\|} dt \\
&= \|x^k - x^*\|^2 \gamma_k \int_0^1 (1-t) dt \\
&= \frac{\gamma_k}{2} \|x^k - x^*\|^2 \\
&\leq \frac{1}{2} \frac{\gamma}{\sqrt{1-2h}} \|x^k - x^*\|^2
\end{aligned}$$

Note:

$$\frac{d}{dt} f(x^* + t(x^k - x^*)) = Df(x^* + t(x^k - x^*))(x^k - x^*)$$

Parameterizing in terms of t is the reason why we need complexity. \square

7.2 Use of Newton in Practice

Remark 7.1. Monotonicity Test (Preliminary)

Newton iterates: $x^{k+1} = x^k - \underbrace{Df(x^k)^{-1} f(x^k)}_{=\Delta x^k}$.

If all goes well: $x^k \rightarrow x^*$, $f(x^*) = 0$.

Check for progress: $\|f(x^{k+1})\| \leq \theta \|f(x^k)\|$ for some $\theta < 1$. Note that this check is not affine-invariant!

Instead:

$$\|Df(x^k)^{-1} f(x^{k+1})\| \leq \theta \underbrace{\|Df(x^k)^{-1} f(x^k)\|}_{=\|\Delta x^k\|}$$

In addition to

$$Df(x^k) \Delta x^k = -f(x^k) \tag{7.2}$$

we also need to solve a second system:

$$Df(x^k) \bar{\Delta} x^{k+1} = -f(x^{k+1})$$

Additional cost: $O(n^2)$ flop if LU factorization is used to solve (7.2).

Remark 7.2. Monotonicity Test (Refined)

- Compute $\bar{\Delta}x^{k+1}$ in addition to Δx^k .
- Check if $\|\bar{\Delta}x^{k+1}\| \leq \theta \|\Delta x^k\|$ for some $\theta < 1$. (Typical value is $\theta = \frac{1}{2}$)
- If not satisfied, use Newton's method with damping

7.3 Newton's Method with Damping

Remark 7.3. Newton's Method: Standard vs. Damped

Standard: $x^{k+1} = x^k + \Delta x^k$

Damped: $x^{k+1} = x^{k+1}(\lambda_k) = x^k + \lambda_k \Delta x^k$, where $0 < \lambda_k \leq 1$ for some damping factor.

What is a suitable strategy for selecting λ_k ?

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8.1 Newton's Method with Damping (Continued)

Remark 8.1. Goal of Damping

$$\|Df(x^k)^{-1}f(x^{k+1})\| \leq \theta \|Df(x^k)^{-1}f(x^k)\| \quad \text{for some } \theta < 1 \quad (8.1)$$

This can always be satisfied if λ_k is chosen small enough.

In the following,

$$\|\mathbf{v}\| := \sqrt{\mathbf{v}^T \mathbf{v}} = \text{Euclidean norm.}$$

Proposition 8.2.

If $f(x^k) \neq \mathbf{0}$, then the Newton increment

$$x^k = Df(x^k)^{-1}f(x^k)$$

is a descent direction for the function

$$\phi(x) := \frac{1}{2} \|Df(x^k)^{-1}f(x)\|^2.$$

at x^k , i.e.

$$\frac{d}{d\lambda} \phi(x^k + \lambda \delta x^k) \Big|_{\lambda=0} < 0.$$

Corollary 8.3.

If $\lambda_k > 0$ is chosen small enough, then

$$x^{k+1} = x^k + \lambda_k \Delta x^k$$

satisfies (8.1). This implies that

$$\phi(x^{k+1}) \leq \theta^2 \phi(x^k).$$

Proof. (of Proposition 8.2)

Note: $f(x^k) \neq 0 \Rightarrow \Delta x^k \neq 0$.

$$\begin{aligned} \frac{d}{d\lambda} \phi(x^k + \lambda \Delta x^k) \Big|_{\lambda=0} &= D_x \phi(x) \Big|_{x=x^k} \Delta x^k \\ D_x \phi(x) &= (Df(x^k)^{-1}f(x))^T Df(x^k)^{-1} Df(x) \\ &= -(\Delta x^k)^T \Delta x^k \\ &= -\|\Delta x^k\|^2 \\ &< 0 \end{aligned}$$

□

Remark 8.4. Practical Strategy for Selecting λ

Determine $0 < \lambda_k \leq 1$ such that the monotonicity test is satisfied with $\theta = 1 - \frac{\lambda_k}{2}$. If possible, use $\lambda_k = 1$ (which implies that $\theta = \frac{1}{2}$).

Procedure

- for $\lambda_k = 1, \frac{1}{2}, \frac{1}{4}, \frac{1}{8}, \dots, \lambda_{\min}$
 - Solve $Df(x^k)\bar{\Delta}x^{k+1}(\lambda_k) = -f(x^k + \lambda_k\Delta x^k)$ and check if

$$\|\bar{\Delta}x^{k+1}(\lambda_k)\| \leq \left(1 - \frac{\lambda_k}{2}\right) \|\Delta x^k\| \quad (8.2)$$

- If (8.2) is satisfied, use λ_k and set $x^{k+1} = x^k + \lambda_k\Delta x^k$
- In the next Newton step, use $\lambda_{k+1} = \min\{1, 2\lambda_k\}$ as the first value for checking (8.2).

8.2 Conditioning

Problem: Input \rightarrow Output (or No Solution)

Abstract Formulation: $X \rightarrow Y$, $x \in X$ is the input, $y = F(x) \in Y$ is the corresponding output.

Example 8.5. General Problem We Are Considering

Solution of $Ay = b$, where $A \in \mathbb{C}^{n \times n}$ is nonsingular, $b \in \mathbb{C}^n$.

$$\begin{aligned} X &= \{x = (A, b) \mid A \in \mathbb{C}^{n \times n} \text{ is nonsingular}, b \in \mathbb{C}^n\} \\ Y &= \mathbb{C}^n \\ F(x) &= F(A, b) := A^{-1}b = y \in \mathbb{C}^n \end{aligned}$$

Definition 8.6. Conditioning of a Problem

The impact of perturbation: $x \rightarrow x + \delta x$, δx is small. In other words, we are solving a slightly wrong problem and we want to know the effect on the output.

Perturbation of input: $\delta x (= x + \delta x - x)$

Perturbation of output: $\delta F(x) (= F(x + \delta x) - F(x))$

The *condition number* of a problem, $\kappa = \kappa(x)$, is a measure of the sensitivity of the output to small changes in the input:

$$\frac{\|\delta F(x)\|}{\|x\|} \leq \kappa(x) \frac{\|\delta x\|}{\|x\|} \quad \text{for small } \|\delta x\|$$

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9.1 Conditioning (Continued)

We have $F : X \rightarrow Y$, $x + \delta x \rightarrow F(x) + \delta F(x)$.

$$\frac{\|\delta F(x)\|}{\|F(x)\|} \leq \kappa(x) \frac{\|\delta x\|}{\|x\|} \quad \text{for small } \|\delta x\|$$

$$\kappa(x) \leq \frac{\|\delta F(x)\|}{\|F(x)\|} \frac{\|x\|}{\|\delta x\|}$$

Here, $\|\cdot\|$ are suitably chosen norms in X and Y .

Definition 9.1. $\kappa(x)$

$$\kappa(x) := \limsup_{\substack{\delta \rightarrow 0 \\ \|\delta x\| \leq \delta}} \frac{\|\delta F(x)\|}{\|F(x)\|} \frac{\|x\|}{\|\delta x\|}$$

$$= \sup_{\delta x} \frac{\|\delta F(x)\|}{\|\delta x\|} \frac{\|x\|}{\|F(x)\|}$$

Definition 9.2. Ill-Conditioned, Well-Conditioned

A problem is said to be *ill-conditioned* if $\kappa(x) \gg 1$ and *well-conditioned* otherwise.

How do we compute the condition number?

Remark 9.3. Conditioning Number for F Differentiable

If $X \subset \mathbb{R}^n$, $Y \subset \mathbb{R}^m$ and $F : X \rightarrow Y$ is differentiable, then

$$\kappa(x) = \|DF(x)\| \frac{\|x\|}{\|F(x)\|}.$$

(Note: $\frac{\|\delta F(x)\|}{\|\delta x\|}$ is kind of like a difference quotient.)

Example 9.4.

Evaluation of $y = F(x) = \sqrt{x}$, $x > 0$.

$$\begin{aligned}F'(x) &= \frac{1}{2} \frac{1}{\sqrt{x}} \\ \kappa(x) &= |F'(x)| \frac{|x|}{|F(x)|} = \frac{1}{2} \cdot \frac{1}{\sqrt{x}} \cdot \frac{x}{\sqrt{x}} \\ &= \frac{1}{2}\end{aligned}$$

So this is well-conditioned for all $x > 0$.

Example 9.5.

Computing $y = F(x) = x_1 - x_2$ for $x = [x_1 \ x_2]^T \in \mathbb{R}^2$.

$F : \mathbb{R}^2 \rightarrow \mathbb{R}$. F is differentiable:

$$DF = [1 \ -1]$$

We compute:

$$\kappa(x) = \|DF(x)\|_2 \frac{\|x\|_2}{|F(x)|} = \sqrt{2} \frac{\sqrt{x_1^2 + x_2^2}}{|x_1 - x_2|}$$

So this problem is ill-conditioned if $x_1 \approx x_2$, $x_1, x_2 \neq 0$.

The reason why $F(x) = x_1 - x_2$ is ill-conditioned is that small perturbations of x_1 and x_2 get amplified:

$$\begin{aligned}F(x + \delta x) &= (x_1 + \delta x_1) - (x_2 + \delta x_2) = \underbrace{x_1 - x_2}_{\approx 0} + \delta x_1 - \delta x_2 \\ &\approx \delta x_1 - \delta x_2\end{aligned}$$

This is called “loss of significant digits.”

Remark 9.6.

Conditioning is only a property of the problem to be solved, not of a specific algorithm for obtaining its solution.

To compute a condition number in Matlab, use `cond(A)`.

9.2 Floating-Point Numbers

Problem: How do we represent $x \in \mathbb{R}$ on a computer?

Finite Storage: irrational numbers like $\sqrt{3} = 1.73205\dots$ or even very large integers cannot be represented exactly!

Base of the representation: $\beta \geq 2$

- $\beta = 2$: binary representation
- $\beta = 10$: decimal representation

If $x \geq 0$ is an integer,

$$\begin{aligned} x &= b_n\beta^n + b_{n-1}\beta^{n-1} + \cdots + b_1\beta + b_0 \\ &=: (b_nb_{n-1}\cdots b_1b_0)_\beta \end{aligned}$$

where the b_i 's are integers such that $0 \leq b_i \leq \beta - 1$.

Example 9.7. 53 in Base 10 and Base 2

Let $x = 53$.

$$\begin{aligned} x &= 5 \cdot 10^1 + 3 \cdot 10^0 = (53)_{10} \\ &= 32 + 16 + 4 + 1 = 1 \cdot 2^5 + 1 \cdot 2^4 + 0 \cdot 2^3 + 1 \cdot 2^2 + 0 \cdot 2^1 + 1 \cdot 2^0 \\ &= (110101)_2 \quad \Rightarrow \text{ 6 bits} \end{aligned}$$

Example 9.8. How many numbers can we store?

Suppose we use 52 bits to store x in a binary representation. How many numbers can we represent exactly?

$$x = (b_{51}b_{50}\cdots b_1b_0)_2$$

Answer: $x = 0, 1, 2, \dots, 2^{52} - 1$.

In general we need to use floating-point representation, even for integers.

9.3 IEEE Floating-Point Standard

Base $\beta = 2$. Any $x \in \mathbb{R}$, $x \neq 0$, can be represented in the binary form

$$x = \pm(1.b_1b_2\cdots) \times 2^p$$

where $b_1, b_2, \dots \in \{0, 1\}$ and p is an integer.

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10.1 Quick Review of the IEEE Standard

$\beta = 2$, $x \in \mathbb{R}$, $x \neq 0$

$$x = (1.b_1 b_2 \dots) \times 2^p$$

where $b_1, b_2, \dots \in \{0, 1\}$ and p is an integer.

Example 10.1. 41.7 in the IEEE standard

$$\begin{array}{ll} \frac{41}{2} = 20 \text{ remainder } 1 & 2 \times 0.7 = 0.4 \text{ remainder } 1 \\ \frac{20}{2} = 10 \text{ remainder } 0 & 2 \times 0.4 = 0.8 \text{ remainder } 0 \\ \frac{10}{2} = 5 \text{ remainder } 0 & 2 \times 0.8 = 0.6 \text{ remainder } 1 \\ \frac{5}{2} = 2 \text{ remainder } 1 & 2 \times 0.6 = 0.2 \text{ remainder } 1 \\ \frac{2}{2} = 1 \text{ remainder } 0 & 2 \times 0.2 = 0.4 \text{ remainder } 0 \\ \frac{1}{2} = 0 \text{ remainder } 1 & \end{array}$$

Thus,

$$\begin{aligned} 41.7 &= (101001.1\overline{0110})_2 \\ &= +(1.01001100110\dots) \times 2^5 \end{aligned}$$

This is its *normalized binary representation*.

10.2 Normalized IEEE Floating-Point Numbers

$$\underbrace{\pm}_{\text{sign}} \underbrace{(1.b_1 b_2 \dots b_N)}_{\text{mantissa}} \times 2^{\overbrace{p}^{\text{exponent}}}$$

where $b_1, b_2, \dots \in \{0, 1\}$ and p can be represented with M bits.

Precision	sign	N (length of mantissa)	M (exponent)	total bits	total bytes
Single	1	23	8	32	4
Double	1	52	11	64	8
Long Double	1	64	15	80	10

From now on we stick with double precision.

A double floating-point number looks like:

$$\pm(1.b_1 b_2 \dots b_{52}) \times 2^p$$

where p is an integer in $[-1022, 1023]$.

Some numbers:

$$1 := +1.\underbrace{0 \ 0 \ \dots \ 0}_{52} \times 2^0$$

$$1 + 2^{-52} := +1.\underbrace{0 \ 0 \ \dots \ 0}_{51} \ 1 \times 2^0$$

10.3 Machine Precision

$1 + 2^{-52}$ is the smallest floating-point number that is still larger than 1. To check this, we loop `for n=0,1,2,...`

:

$$1 + 2^{-n} \stackrel{?}{>} 1$$

It should stop for $n = 53$.

Definition 10.2. *Machine Precision*

Machine precision (or *machine epsilon*) is a measure of precision. Its value is

$$\epsilon_{\text{mach}} := 2^{-52} = (1 + 2^{-52} - 1)$$

(`eps` in Matlab)

Note:

$$\epsilon_{\text{mach}} = 2.220 \dots \times 10^{-16}$$

which is why double precision has ≈ 16 significant digits.

10.4 Floating-Point Representation

$$x \in \mathbb{R} \xrightarrow{\text{round to nearest}} fl(x) = \text{floating-point number closest to } x$$

Remark 10.3. *Chopping vs. Rounding to Nearest*

Consider

$$x = \pm 1.b_1 b_2 \dots b_{52} b_{53} \dots \times 2^p$$

- Chopping:
 - $x \rightarrow \pm 1.b_1 b_2 \dots b_{52} \times 2^p$
- Round down: same as chopping
- Round up: 1 is added to b_{52}
- Rounding to Nearest:
 - If $b_{53} = 0$, round down
 - If $b_{53} = 1$ and $b_j = 1$ for at least one of $j > 53$, round up
 - If $b_{53} = 1$ and $b_j = 0$ for all $j > 53$:
 - * round up if $b_{52} = 1$
 - * round down if $b_{52} = 0$

Definition 10.4. Relative Error

For $x \neq 0$, the *relative error* of $fl(x)$ is

$$\left| \frac{fl(x) - x}{x} \right| \leq \frac{1}{2} \epsilon_{\text{mach}} = 2^{-53} = 1.11\dots \times 10^{-16}$$

Example 10.5. Representation of 41.7 with Rounding

$$41.7 = +1.\underbrace{010011}_{6} \underbrace{0110}_{4} \underbrace{0110}_{4} \dots \underbrace{0110}_{4} \underbrace{01}_{b_{52}} | 10 0110 \dots \times 2^5$$
$$fl(41.7) = +1.010011 0110 0110 \dots 0110 01 \times 2^5$$

11 10-17-11

11.1 Representation of the Exponent p

Double precision uses 11 bits:

$$e_{10}2^{10} + e_92^9 + \cdots + e_12^1 + e_02^0, \quad e_0, e_1, \dots, e_{10} \in \{0, 1\}$$

We can represent all integers from 0 to $2^{11} - 1 = 2047$. BUT p is written in the shifted form:

$$p = \underbrace{-1023}_{\text{fixed}} + e_{10}2^{10} + e_92^9 + \cdots + e_12^1 + e_02^0$$

Thus, p ranges from -1023 to 1024 .

11.2 Machine Representation

s	e_0	e_1	\cdots	e_{10}	$ $	b_1	b_2	\cdots	b_{52}
-----	-------	-------	----------	----------	-----	-------	-------	----------	----------

$$\begin{aligned} x \neq \pm 0, \pm \infty : & \quad s \ e_0 \ e_1 \ \dots \ e_{10} \mid b_1 \ b_2 \ \dots \ b_{52} \\ x = \pm 0 : & \quad s \ 0 \ 0 \ \dots \ 0 \mid 0 \ 0 \ \dots \ 0 \\ x = \pm \infty : & \quad s \ 1 \ 1 \ \dots \ 1 \mid 0 \ 0 \ \dots \ 0 \end{aligned}$$

Undefined Cases:

$$\begin{aligned} \text{NaN: } & s \ 1 \ 1 \ \dots \ 1 \mid b_1 \ b_2 \ \dots \ b_{52}, \quad \text{at least one } b_j = 1 \\ \text{subnormal floating point numbers: } & \underbrace{s \ 0 \ 0 \ \dots \ 0}_{< 2^{-1022}} \mid b_1 \ b_2 \ \dots \ b_{52}, \quad \text{at least one } b_j = 1 \\ & 2^{-1022} : \ 1 \ 1 \ 0 \ \dots \ 0 \mid 0 \ 0 \ \dots \ 0 \end{aligned}$$

11.3 Floating Point Arithmetic

Example 11.1. $41.7 \oplus 10.425$

$$\begin{aligned} x &= fl(41.7) = 1.010011 \ 0110 \ 0110 \ \dots \ 0110 \ 01 \times 2^5 \\ y &= fl(\underbrace{10.425}_{=41.7/4}) = 1.010011 \ 0110 \ 0110 \ \dots \ 0110 \ 01 \times 2^3 \end{aligned}$$

Addition of x and y in floating-point arithmetic:

$$\begin{aligned} x &= +1.0100110110 \dots 011001 \mid \times 2^5 \\ y &= +0.010100110110 \dots 0110 \mid 01 \times 2^5 \\ x + y &= 1.101000010 \dots 00000 \mid 01 \times 2^5 \\ \underbrace{fl(x + y)}_{=x \oplus y \neq x + y} &= 1.101000010 \dots 00000 \mid \times 2^5 \end{aligned}$$

<u>Exact Arithmetic</u>	<u>floating-point arithmetic</u>
$x + y$	$x \oplus y$
$x - y$	$x \ominus y$
$x \times y$	$x \otimes y$
x/y	$x \oslash y$

For all floating-point numbers x, y :

$$\begin{aligned} x \oplus y &= (x + y)(1 + \epsilon) && \text{where } |\epsilon| \leq \frac{1}{2}\epsilon_{\text{mach}} = 2^{-53} \\ x \ominus y &= (x - y)(1 + \epsilon) && \text{where } |\epsilon| \leq \frac{1}{2}\epsilon_{\text{mach}} \\ x \otimes y &= (x \times y)(1 + \epsilon) && \text{where } |\epsilon| \leq \frac{1}{2}\epsilon_{\text{mach}} \\ (y \neq 0) \quad x \oslash y &= (x/y)(1 + \epsilon) && \text{where } |\epsilon| \leq \frac{1}{2}\epsilon_{\text{mach}} \end{aligned}$$

Remark 11.2.

The relative error of the floating-point implementations of the four basic arithmetic operations is bounded by

$$\frac{1}{2}\epsilon_{\text{mach}} = 2^{-53} = 1.11 \dots \times 10^{-16}.$$

11.4 Catastrophic Effects of Round-Off Errors

Loss of Significant Digits

Occurs in subtraction of nearly equal numbers, i.e.

$$x - y, \quad \text{where } x \approx y.$$

Example 11.3. Relative Error of $x - y$ when $x \approx y$

$$\begin{aligned} x &= 1 + 2^{-52} + 2^{-53} + 2^{-54} &= +1.0 \dots 01 \mid 110 \dots \times 2^0 \\ y &= 1 + 2^{-54} &= +1.0 \dots 00 \mid 010 \dots \times 2^0 \end{aligned}$$

Exact arithmetic:

$$x - y = 2^{-52} + 2^{-53} = 3 \times 2^{-53}$$

Floating-point arithmetic:

$$\begin{aligned} fl(x) &= +1.0 \dots 010 \mid \times 2^0 \\ fl(y) &= +1.0 \dots 000 \mid \times 2^0 \\ fl(x) - fl(y) &= +0.0 \dots 010 \mid \times 2^0 = +1.0 \dots 0 \times 2^{-51} \end{aligned}$$

Relative error:

$$\left| \frac{2^{-51} - 3 \times 2^{-53}}{3 \times 2^{-53}} \right| = \frac{1}{3}$$

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12.1 Loss of Significant Digits

$$x - y, \quad x \approx y$$

In practice, this problem can often be avoided easily.

Example 12.1. *Trouble Is Avoidable*

Evaluation of

$$f(x) = \frac{1 - \cos x}{\sin^2 x} \quad \text{for } 0 < x < \pi$$

Problem: $\cos x \approx 1$ for $x \approx 0$.

$$\begin{aligned} f(x) &= \frac{1 - \cos x}{\sin^2 x} \cdot \frac{1 + \cos x}{1 + \cos x} = \frac{1 - \cos^2 x}{\sin^2 x} \cdot \frac{1}{1 + \cos x} \\ &= \frac{1}{1 + \cos x} \end{aligned}$$

But $\cos x \approx -1$ for $x \approx \pi$.

⇒ use

$$f(x) = \begin{cases} \frac{1}{1+\cos x} & 0 < x \leq \frac{\pi}{2} \\ \frac{1-\cos x}{\sin^2 x} & \frac{\pi}{2} < x < \pi \end{cases}$$

Example 12.2. Huge Intermediate Quantities

$$f(x) = \frac{1}{x} \left(\frac{\frac{1}{x}}{\frac{1}{x} - 1} - 1 \right), \quad 0 < x < 1$$

Suppose we want to evaluate this at $x = 2^{-54}$.

In exact arithmetic:

$$\begin{aligned} y = f(2^{-54}) &= 2^{54} \left(\frac{2^{54}}{2^{54} - 1} - 1 \right) \\ &= \frac{1}{1 - 2^{-54}} \\ &\approx 1 \end{aligned}$$

In floating-point arithmetic:

$$\begin{aligned} \tilde{y} &= 2^{54} \otimes ((2^{54} \oslash \underbrace{(2^{54} \ominus 1)}_{= 0}) \ominus 1) \\ &= 2^{54} \otimes ((\underbrace{2^{54} \oslash 2^{54}}_{= 1}) \ominus 1) \\ &= 2^{54} \otimes (\underbrace{1 \ominus 1}_{= 0}) \\ &= 2^{54} \otimes 0 \\ &= 0 \end{aligned}$$

$$\text{Relative Error} = \left| \frac{\tilde{y} - y}{y} \right| = 1$$

The problem is the huge intermediate quantity 2^{54} .

Remedy:

$$f(x) = \frac{1}{x} \left(\frac{1}{\frac{1}{x} - 1} \right) = \frac{1}{1 - x}$$

In floating-point arithmetic:

$$\begin{aligned} \tilde{y} &= 1 \oslash (1 \ominus 2^{-54}) = 1 \oslash 1 \\ &= 1 \end{aligned}$$

$$\text{Relative Error} = \left| \frac{\tilde{y} - y}{y} \right| = \left| \frac{1 - \frac{1}{1-2^{-54}}}{\frac{1}{1-2^{-54}}} \right| = 2^{-54} < \epsilon_{\text{mach}}$$

12.2 Stability

Abstract formulation of a problem:

$$\text{Input } x \in X \xrightarrow[F:X \rightarrow Y]{F} F(x) = y \in Y, \quad y \text{ is the output}$$

An algorithm in floating-point arithmetic:

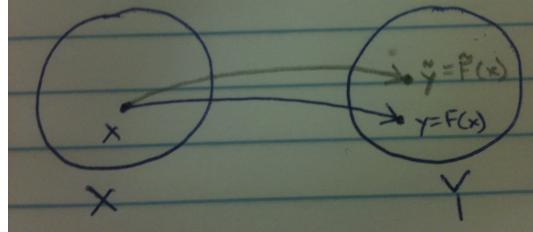
$$\tilde{F} : X \rightarrow Y$$

\tilde{F} captures all errors:

- $x \rightarrow fl(x)$
- round-off errors of floating-point arithmetic
- possible approximation error (e.g. stopping Newton's method)

Exact solution of the problem: $y = F(x)$

Computed solution: $\tilde{y} = \tilde{F}(x)$



Relative error of the computed solution:

$$\frac{\|\tilde{y} - y\|}{\|y\|} = \frac{\|\tilde{F}(x) - F(x)\|}{\|F(x)\|}$$

Ideally we would like to have

$$\frac{\|\tilde{F}(x) - F(x)\|}{\|F(x)\|} = O(\epsilon_{\text{mach}}) \quad \text{for all } x \in X \quad (12.1)$$

Recall: $g(\epsilon) = O(\epsilon) \Leftrightarrow$ there exists a constant C such that $\|g(\epsilon)\| \leq C|\epsilon|$ for $\epsilon \rightarrow 0$.

If the condition number is 1, then (12.1) is possible. If the condition number is big, i.e. the problem is ill-conditioned, then (12.1) is unrealistic.

Instead...

Definition 12.3. *Stable*

The “algorithm” $\tilde{F} : X \rightarrow Y$ for solving the problem is *stable* if for each $x \in X$, there exists an $\tilde{x} \in X$ such that

$$(1) \quad \frac{\|\tilde{F}(x) - F(\tilde{x})\|}{\|F(\tilde{x})\|} = O(\epsilon_{\text{mach}})$$

$$(2) \quad \frac{\|\tilde{x} - x\|}{\|x\|} = O(\epsilon_{\text{mach}})$$

Meaning of Stable

A stable algorithm gives nearly the right solution to nearly the right problem.

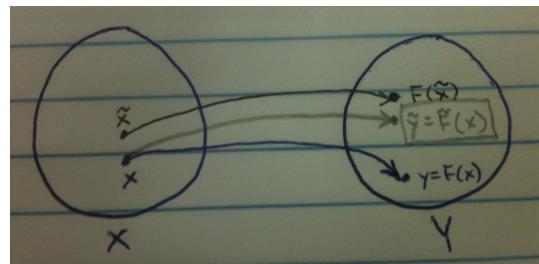


Figure 2: $\tilde{y} = \tilde{F}(x)$ is the computed solution.

This is a realistic goal.

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13.1 Backward Stability

Definition 13.1. Backward Stability

A stronger form of stability:

The “algorithm” $\tilde{F} : X \rightarrow Y$ for solving the problem $F : X \rightarrow Y$ is said to be *backward stable* if for each input there exists an $\tilde{x} \in X$ such that

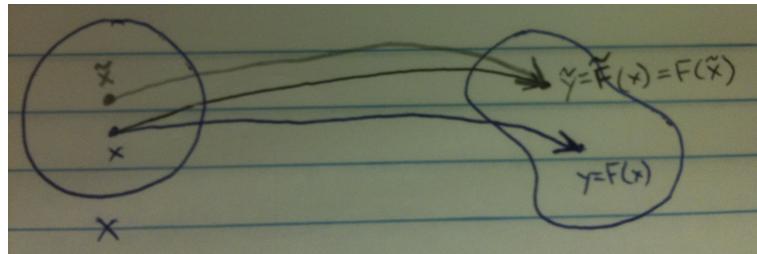
$$\tilde{F}(x) = F(\tilde{x})$$

and

$$\frac{\|\tilde{x} - x\|}{\|x\|} = O(\epsilon_{\text{mach}})$$

Meaning of Backward Stability:

A backward stable algorithm gives the right solution to nearly the right problem.



Example 13.2. Backward Stability of Floating-Point Subtraction

$$y = F(x) = x_1 - x_2, \quad x = \begin{bmatrix} x_1 \\ x_2 \end{bmatrix} \in \mathbb{R}^2$$

$$\tilde{y} = \tilde{F}(x) = fl(x_1) \ominus fl(x_2)$$

$$= [x_1(1 + \epsilon_1) - x_2(1 + \epsilon_2)](1 + \epsilon_3) \quad \text{where } |\epsilon_1|, |\epsilon_2|, |\epsilon_3| \leq \frac{1}{2}\epsilon_{\text{mach}}$$

$$= \underbrace{x_1(1 + \epsilon_1)(1 + \epsilon_3)}_{=: \tilde{x}_1} - \underbrace{x_2(1 + \epsilon_2)(1 + \epsilon_3)}_{=: \tilde{x}_2}$$

$$= F(\tilde{x})$$

$$\left| \frac{\tilde{x}_1 - x_1}{x_1} \right| = |(1 + \epsilon_1)(1 + \epsilon_3) - 1| = |\epsilon_1 + \epsilon_3 + \epsilon_1 \epsilon_3|$$

$$\leq |\epsilon_1| + |\epsilon_3| + |\epsilon_1||\epsilon_3| = O(\epsilon_{\text{mach}})$$

$$\left| \frac{\tilde{x}_2 - x_2}{x_2} \right| = O(\epsilon_{\text{mach}})$$

by similar analysis

The floating point algorithm for subtraction is backward stable.

Remark 13.3. Backward Stability of Floating-Point Addition, Multiplication, and Division

Similar to Example 13.2, we can show that addition, multiplication, and division are backward stable.

Example 13.4. Stable, but not Backward Stable

$$y = F(x) = x + 1, \quad x \in \mathbb{R}$$

$$\begin{aligned} \tilde{y} = \tilde{F}(x) &= fl(x) \oplus 1 = [x(1 + \epsilon_1) + 1](1 + \epsilon_2) && \text{where } |\epsilon_1|, |\epsilon_2| \leq \frac{1}{2}\epsilon_{\text{mach}} \\ &= \underbrace{x(1 + \epsilon_1)(1 + \epsilon_2) + \epsilon_2}_{{=: \tilde{x}}} + 1 = \tilde{x} + 1 = F(\tilde{x}) \left| \frac{\tilde{x} - x}{x} \right| && = \left| \frac{x(1 + \epsilon_1)(1 + \epsilon_2) + \epsilon_2 - x}{x} \right| \\ &= \left| \epsilon_1 + \epsilon_2 + \epsilon_1 \epsilon_2 + \frac{\epsilon_2}{x} \right| = O(\epsilon_{\text{mach}}) + O\left(\frac{\epsilon_{\text{mach}}}{|x|}\right) \\ &\neq O(\epsilon_{\text{mach}}) \quad \text{if } x \approx 0 \end{aligned}$$

Thus, this algorithm is not backward stable. Is it stable?

$$\begin{aligned} F(x) &= x + 1 \\ \tilde{F}(x) &= \underbrace{x(1 + \epsilon_1)(1 + \epsilon_2)}_{{=: \tilde{x}}} + \epsilon_2 + 1 \\ \tilde{F}(\tilde{x}) - F(\tilde{x}) &= (\tilde{x} + 1)(1 + \epsilon_2) - (\tilde{x} + 1) = (\tilde{x} + 1)\epsilon_2 \\ &= F(\tilde{x})\epsilon_2 \\ |\tilde{F}(x) - F(\tilde{x})| &= |F(\tilde{x})||\epsilon_2| = |F(\tilde{x})|O(\epsilon_{\text{mach}}) \\ \tilde{x} - x &= x\epsilon_1 \\ |\tilde{x} - x| &= |x|O(\epsilon_{\text{mach}}) \end{aligned}$$

Thus, it is stable.

13.2 Accuracy of Backward Stable Algorithms

Theorem 13.5.

Let the problem $F : X \rightarrow Y$ have condition number $\kappa = \kappa(x)$, $x \in X$, and let $\tilde{F} : X \rightarrow Y$ be a backward stable algorithm for solving the problem. Then:

$$\underbrace{\frac{\|\tilde{F}(x) - F(x)\|}{\|F(x)\|}}_{=\frac{\|\tilde{y}-y\|}{\|y\|}} = O(\kappa(x)\epsilon_{\text{mach}}) \quad \text{for all } x \in X.$$

Thus, if the condition number is really high, then even a backward stable algorithm cannot save you.

14 10-24-11

14.1 Backward Stability (Continued)

Theorem 14.1.

$$F : X \rightarrow Y, \quad \kappa(x) \\ \tilde{F} : X \rightarrow Y \quad \text{backward stable}$$

$$\frac{\|\tilde{F}(x) - F(x)\|}{\|F(x)\|} = O(\kappa(x)\epsilon_{\text{mach}}) \quad \forall x \in X$$

Proof. Let $x \in X$ and let \tilde{F} be backward stable. Then

$$\tilde{F}(x) = F(\tilde{x}) \quad \text{for some } \tilde{x} \in X \text{ with } \frac{\|\tilde{x} - x\|}{\|x\|} = O(\epsilon_{\text{mach}}).$$

Recall:

$$\kappa(x) = \sup_{\delta x} \frac{\|\delta F(x)\|}{\|\delta x\|} \cdot \frac{\|x\|}{\|F(x)\|}$$

Then

$$\begin{aligned} \frac{\|\tilde{F}(x) - F(x)\|}{\|F(x)\|} &= \frac{\|F(\tilde{x}) - F(x)\|}{\|F(x)\|} && \text{(backward stability)} \\ &= \frac{\|\delta F(x)\|}{\|F(x)\|} && (\delta x := \tilde{x} - x) \\ &= \left(\frac{\|\delta F(x)\|}{\|\delta x\|} \cdot \frac{\|x\|}{\|F(x)\|} \right) \frac{\|\tilde{x} - x\|}{\|x\|} \\ &\leq (\kappa(x) + \underbrace{o(1)}_{\rightarrow 0 \text{ as } \epsilon_{\text{mach}} \rightarrow 0}) \underbrace{\frac{\|\tilde{x} - x\|}{\|x\|}}_{=O(\epsilon_{\text{mach}})} \\ &= O(\kappa(x)\epsilon_{\text{mach}}) \end{aligned}$$

□

14.2 Norms

Definition 14.2. Norm

A *norm* on \mathbb{C}^n is a function

$$\|\cdot\| : \mathbb{C}^n \rightarrow \mathbb{R}$$

such that for all $x, y \in \mathbb{C}^n$, $\alpha \in \mathbb{C}$:

1. $\|x\| \geq 0$, and $\|x\| = 0$ if and only if $x = 0$
2. $\|x + y\| \leq \|x\| + \|y\|$
3. $\|\alpha x\| = |\alpha| \|x\|$

Example 14.3. Some Norms

$$\begin{aligned}\|x\|_1 &= \sum_{i=1}^n |x_i| \\ \|x\|_2 &= \sqrt{\sum_{i=1}^n |x_i|^2} \\ \|x\|_\infty &= \max_{1 \leq i \leq n} |x_i|\end{aligned}$$

Example 14.4. Vec Norm

Let $A \in \mathbb{C}^{n \times n}$. $\text{vec}(A) := [a_{11} \ a_{21} \ \cdots \ a_{n1} \ a_{12} \ \cdots \ a_{nn}]^T \in \mathbb{C}^{n^2}$. $\|A\| := \|\text{vec}(A)\|$, where $\|\cdot\| : \mathbb{C}^{n^2} \rightarrow \mathbb{R}$.

Example 14.5. Frobenius Norm

$$\|A\|_F = \|\text{vec}(A)\|_2 = \sqrt{\sum_{j,k=1}^n |a_{jk}|^2}$$

14.3 Matrix Norms

Definition 14.6. Matrix Norm

Any *matrix norm* satisfies:

1. $\|A\| \geq 0$, with $\|A\| = 0$ if and only if $A = \mathbf{0}$
2. $\|A + B\| \leq \|A\| + \|B\|$
3. $\|\alpha A\| = |\alpha| \|A\|$

Definition 14.7. Induced Matrix Norm

Often more useful are *induced matrix norms*:

Let $\|\cdot\| : \mathbb{C}^n \rightarrow \mathbb{R}$ be a matrix norm on \mathbb{C}^n . Set

$$\text{lub}(A) = \text{lub}_{\|\cdot\|}(A) := \max_{x \neq 0} \frac{\|Ax\|}{\|x\|} = \max_{\|x\|=1} \|Ax\|$$

It is easy to see that $\text{lub}(A)$ is indeed a matrix norm.

Note: $\|Ax\| \leq \text{lub}_{\|\cdot\|}(A) \|x\|$ for all $A \in \mathbb{C}^{n \times n}$, $x \in \mathbb{R}^n$.

Norm on \mathbb{C}^n Induced Matrix Norm

$\ \cdot\ _1$	$\ A\ _1 := \text{lub}_{\ \cdot\ }(A) = \max_{1 \leq k \leq n} \sum_{j=1}^n a_{jk} $	(max column sum)
$\ \cdot\ _\infty$	$\ A\ _\infty := \text{lub}_{\ \cdot\ _\infty}(A) = \max_{1 \leq j \leq n} \sum_{k=1}^n a_{jk} $	(max row sum)
$\ \cdot\ _2$	$\ A\ _2 := \text{lub}_{\ \cdot\ }(A) = \max_{\ x\ _2=1} \ Ax\ _2 = \sigma_{\max}(A) = \text{the largest singular value of } A$	

$\sigma_{\max} = \sqrt{\lambda_{\max}(A^H A)}$

Theorem 14.8. lub Norms Are Submultiplicative

lub norms are *submultiplicative*:

$$\text{lub}(AB) \leq \text{lub}(A)\text{lub}(B) \quad \forall A, B \in \mathbb{C}^{n \times n}$$

Proof.

$$\text{lub}(AB) = \max_{x \neq 0} \frac{\|ABx\|}{\|x\|}$$

For $x \neq 0$, $y := Bx \neq 0$

$$\begin{aligned} \frac{\|ABx\|}{\|x\|} &= \frac{\|ABx\|}{\|Bx\|} \cdot \frac{\|Bx\|}{\|x\|} = \frac{\|Ay\|}{\|y\|} \frac{\|Bx\|}{\|x\|} \\ &\leq \text{lub}(A)\text{lub}(B) \end{aligned}$$

□

15 10-26-11

15.1 Conditioning of $Ax = b$

Let $A \in \mathbb{C}^{n \times n}$ be nonsingular.

Let $b \in \mathbb{C}^n$ be fixed, and consider small perturbations δA of A .

$$\begin{aligned} \left. \begin{aligned} Ax &= b \\ (A + \delta A)(x + \delta x) &= b \end{aligned} \right\} &\Rightarrow A(\delta x) + (\delta A)x + \underbrace{(\delta A)(\delta x)}_{\approx 0} = 0 \\ A(\delta x) &\approx -(\delta A)x \quad \Rightarrow \quad \delta x \approx -A^{-1}(\delta A)x \end{aligned}$$

Let $\|\cdot\|$ be any vector norm on \mathbb{C}^n and $\|\cdot\|$ be the associated lub norm on $\mathbb{C}^{n \times n}$.

$$\begin{aligned} \|\delta x\| &\approx \|A^{-1}(\delta A)x\| \leq \|A^{-1}(\delta A)\| \|x\| \\ &\leq \|A^{-1}\| \|\delta A\| \|x\| \\ \frac{\|\delta x\|}{\|x\|} \frac{\|A\|}{\|\delta A\|} &\lesssim \|A\| \|A^{-1}\| =: \kappa(A) (= \kappa_{\|\cdot\|}(A)) \end{aligned}$$

Recall our definition of $\kappa(x)$:

$$\sup_{\delta A} \frac{\|\delta x\|}{\|x\|} \frac{\|A\|}{\|\delta A\|} = \kappa(A) = \|A\| \|A^{-1}\|$$

Theorem 15.1.

Let $b \in \mathbb{C}^n$ be fixed. Let $\|\cdot\|$ be a vector norm on \mathbb{C}^n and let $\|\cdot\|$ be the associated lub norm. Let A be nonsingular. The conditioning number of solving $Ax = b$ is given by

$$\kappa(A) = \|A\| \|A^{-1}\|.$$

Corollary 15.2.

$\kappa(A) \geq 1$ for all A (since we are using lub norms).

Proof.

$$\begin{aligned} I &= AA^{-1} \\ \|I\| &= \|AA^{-1}\| \\ &\leq \|A\| \|A^{-1}\| \\ &= \kappa(A) \\ \|I\| &= \max_{\|x\|=1} \|Ix\| = 1 \end{aligned}$$

□

15.2 Stability of LU Factorization

Let $A \in \mathbb{C}^{n \times n}$ be nonsingular. Assume that A has an LU factorization: $A = LU$.

Then: for all “sufficiently small” ϵ_{mach} , LU factorization without partial pivoting in floating-point arithmetic encounters no zero pivots. The algorithm will run to completion, computing a lower-triangular matrix \tilde{L} and an upper-triangular matrix \tilde{U} such that

$$\tilde{L}\tilde{U} = A + \delta A \quad \text{and} \quad \frac{\|\delta A\|}{\|L\| \cdot \|U\|} = O(\epsilon_{\text{mach}}).$$

If $\|L\| \cdot \|U\| = O(\|A\|)$, then

$$\frac{\|\delta A\|}{\|A\|} = O(\epsilon_{\text{mach}})$$

and LU factorization without pivoting is backward stable.

However, for general nonsingular $A \in \mathbb{C}^{n \times n}$,

$$\|L\| \cdot \|U\| \neq O(\|A\|)$$

In fact, without pivoting $\|L\|$ can be arbitrarily large (see the 2×2 example we did).

\Rightarrow LU factorization without pivoting is not backward stable (and not even stable).

15.3 LU Factorization with Partial Pivoting

Partial pivoting ensures that $|l_{jk}| \leq 1$ for all j, k $\Rightarrow \|L\| = O(1)$. But what can we say about $\|U\|$?

Define the so-called growth factor.

Definition 15.3. Growth Factor

$$\rho = \rho(A) = \frac{\max_{j,k=1,2,\dots,n} |u_{jk}|}{\max_{j,k=1,2,\dots,n} |a_{jk}|}$$

Theorem 15.4.

Let $A \in \mathbb{C}^{n \times n}$ be a nonsingular matrix. Then LU factorization with partial pivoting in floating-point arithmetic generates matrices \tilde{L} , \tilde{U} , and \tilde{P} such that

$$\tilde{L}\tilde{U} = \tilde{P}A + \delta A$$

and $\frac{\|\delta A\|}{\|A\|} = O(\rho\epsilon_{\text{mach}})$, where $\rho = \rho(A)$ is the growth factor.

Consequence:

If $\rho(A)$ is bounded by a constant for all nonsingular matrices $A \in \mathbb{C}^{n \times n}$, then

$$\frac{\|\delta A\|}{\|A\|} = O(\epsilon_{\text{mach}})$$

and thus, LU factorization with partial pivoting is backward stable.

This is indeed the case. It can be shown that

$$\rho(A) \leq 2^{n-1} \quad \text{for all nonsingular matrices } A \in \mathbb{C}^{n \times n}$$

and the inequality is sharp. However, in practice, $\rho(A)$ does not approach this bound.

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16.1 Backward Stability of LU Factorization with Partial Pivoting (Continued)

Example 16.1. Large $\rho(A)$

$$A = \begin{bmatrix} 1 & 0 & 0 & \cdots & 0 & 1 \\ -1 & 1 & 0 & \cdots & 0 & 1 \\ -1 & -1 & 1 & \ddots & 0 & 1 \\ \vdots & \vdots & \ddots & \ddots & \vdots & \vdots \\ \vdots & \vdots & \vdots & \ddots & \ddots & \vdots \\ -1 & -1 & -1 & \cdots & -1 & 1 \end{bmatrix} \in \mathbb{R}^{n \times n}$$

$$A \rightarrow \begin{bmatrix} 1 & 0 & 0 & \cdots & 0 & 1 \\ 0 & 1 & 0 & \cdots & 0 & 2 \\ 0 & -1 & 1 & \ddots & \vdots & \vdots \\ \vdots & \vdots & \vdots & \ddots & 0 & 2 \\ 0 & -1 & -1 & \ddots & 1 & 2 \\ 0 & -1 & -1 & \cdots & -1 & 2 \end{bmatrix}$$

$$U = \begin{bmatrix} 1 & 0 & 0 & \cdots & 0 & 1 \\ 0 & 1 & 0 & \cdots & 0 & 2 \\ 0 & 0 & 1 & \ddots & \vdots & \vdots \\ \vdots & \vdots & \ddots & \ddots & 0 & 2^{n-3} \\ 0 & 0 & 0 & \ddots & 1 & 2^{n-2} \\ 0 & 0 & 0 & \cdots & 0 & 2^{n-1} \end{bmatrix}$$

$$L = \begin{bmatrix} 1 & 0 & 0 & \cdots & 0 & 0 \\ -1 & 1 & 0 & \cdots & 0 & 0 \\ -1 & -1 & 1 & \ddots & \vdots & \vdots \\ \vdots & \vdots & \ddots & \ddots & 0 & 0 \\ -1 & -1 & -1 & \ddots & 1 & 0 \\ -1 & -1 & -1 & \cdots & -1 & 1 \end{bmatrix}$$

$$\rho(A) = 2^{n-1}$$

Example 16.2. Large $\rho(A)$ (Continued)

For nonsingular matrices $A \in \mathbb{C}^{n \times n}$ of all sizes $n \geq 1$, LU factorization with partial pivoting is not backward stable!

Reason: $\frac{\|\delta A\|}{\|A\|} = O(\rho(A)\epsilon_{\text{mach}}) \neq O(\epsilon_{\text{mach}})$, since $\rho(A) = 2^{n-1} \rightarrow \infty$ as $n \rightarrow \infty$.

Good News: In practice, matrices with large $\rho(A)$ have never occurred!

Remark 16.3. Backward Stability of LU Factorization

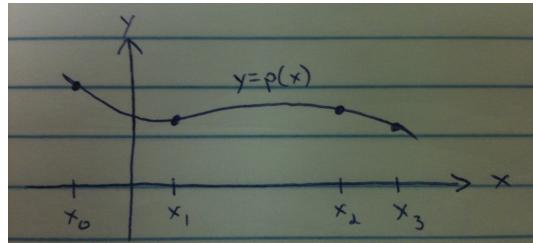
LU factorization with partial pivoting is backward stable in practice!

16.2 Interpolation

Problem: Given $n + 1$ data points

$$(x_j, y_j), \quad j = 0, 1, \dots, n \quad (16.1)$$

where $x_j, y_j \in \mathbb{R}$ and $x_0 < x_1 < x_2 < \dots < x_n$.



Find a “simple” function p that interpolates the data.

16.2.1 Polynomial Interpolation

Remark 16.4. Polynomial Notation

$$P_n \in \prod_n := \left\{ P(x) = c_0 + c_1 x + \dots + c_n x^n \mid c_0, c_1, \dots, c_n \in \mathbb{R} \right\}$$

= set of all (real) polynomials of degree $\leq n$

Theorem 16.5.

There is a unique $P_n \in \prod_n$ that interpolates the data (16.1):

$$P_n(x) = \sum_{j=0}^n y_j \prod_{\substack{k=0 \\ k \neq j}}^n \frac{x - x_k}{x_j - x_k}$$

(Lagrange interpolation formula)

Proof.

$$\begin{aligned} P_n(x_l) &= \sum_{j=0}^n y_j \prod_{\substack{k=0 \\ k \neq j}}^n \frac{x_l - x_k}{x_j - x_k} \\ \prod_{\substack{k=0 \\ k \neq j}}^n \frac{x_l - x_k}{x_j - x_k} &= \begin{cases} 0 & l \neq j \\ 1 & l = j \end{cases} \\ P_n(x_l) &= y_l \quad \forall l \end{aligned}$$

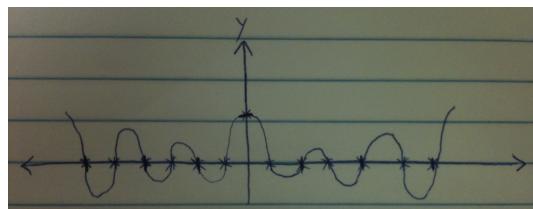
P_n is unique: Let $Q \in \prod_n$ be another interpolating polynomial. Then

$$\begin{aligned} D(x) &= P_n(x) - Q(x) \in \prod_n \\ D(x_j) &= P_n(x_j) - Q(x_j) = y_j - y_j = 0, \quad j = 0, 1, \dots, n \end{aligned}$$

Thus, D has at least $n+1$ zeros. Since D is a polynomial of degree n , this implies that $D = 0 \Rightarrow P_n = Q$. \square

But...

Polynomials are usually not flexible enough.



Remedy: piecewise polynomials \Rightarrow *splines*

16.2.2 Splines

Given: data points $(x_0, y_0), (x_1, y_1), (x_2, y_2), \dots, (x_n, y_n)$, where $a := x_0 < x_1 < x_2 < \dots < x_n =: b$ and $n \geq 1$.

Definition 16.6. Spline

A function $S : [a, b] \rightarrow \mathbb{R}$ is called a *spline* of degree $k-1$ (order k) if

- (a) $S \in C^{k-2}[a, b]$
- (b) On each interval $[x_{j-1}, x_j]$, $j = 1, 2, \dots, n$, $S(x) \in \prod_{k-1}$

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$$a = x_0 < x_1 < x_2 < \dots < x_n = b$$

$$S \in C^{k-2}[a, b]$$

$$S(x) \in \prod_{k=1}^n, \quad x_{j-1} \leq x_j$$

$k \geq 2$, k order, $k - 1$ degree

17.1 Working with Splines

Example 17.1. Linear & Cubic Splines

$k = 2 \Rightarrow$ linear spline

$k = 4, \Rightarrow$ cubic spline

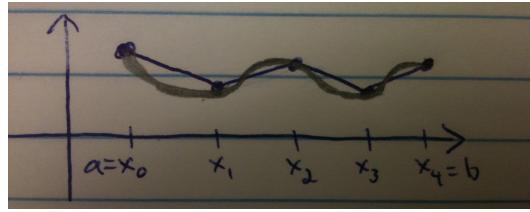


Figure 3: Linear spline (pen) and cubic spline (pencil).

Remark 17.2. Counting Degrees of Freedom

On each $[x_{j-1}, x_j] : S = S_j \in \prod_{k=1}^n \Rightarrow k$ coefficients.

\Rightarrow The total number of parameters = nk coefficients.

$S \in C^{k-2}$. For each x_j , $j = 1, 2, \dots, n - 1$, $S_j^{(i)}(x_j) = S^{(i)}(x_j - 0) = S^{(i)}(x_j + 0) = S_{j+1}^{(i)}(x_j)$, $i = 0, 1, \dots, k - 2$. \Rightarrow This is a total of $(n - 1)(k - 1)$ conditions.

Use a spline of degree $k - 1$ to interpolate the data (x_0, y_0) , (x_1, y_1) , \dots , (x_n, y_n) :

$$S(x_j) = y_j, \quad j = 0, 1, \dots, n$$

\Rightarrow Another $n + 1$ conditions.

$$\begin{aligned} \# \text{ coefficients} - \# \text{ conditions} &= nk - (n - 1)(k - 1) - (n - 1) \\ &= k - 2 \end{aligned}$$

Thus, we need to impose $k - 2$ more conditions in order to get a unique spline.

Example 17.3.

For $k = 2$, there is a unique linear spline that interpolates the data.

For $k = 4$, we need 2 additional conditions to have a unique cubic spline that interpolates the data.

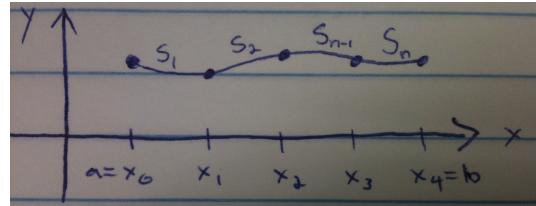
17.1.1 Cubic Splines

$k = 4$, $S \in C^2[a, b]$,

$$S(x) = S_j(x) \in \prod_3 \text{ on } [x_{j-1}, x_j], \quad j = 1, 2, \dots, n.$$

Remark 17.4. Additional 2 Conditions

1. Natural spline: $S''(a) = S_1''(a) = 0$, $S''(b) = S_n''(b) = 0$.
2. Not-a-knot condition: we need that $n \geq 3$. $S_1'''(x_1) = S_2'''(x_1)$, $S_{n-1}'''(x_{n-1}) = S_n'''(x_{n-1})$.



Recall: $S_1, S_2 \in \prod_3$, $S_1, S_2 \in C^2$, $S_1^{(i)}(x_1) = S_2^{(i)}(x_2)$ for $i = 0, 1, 2, 3$. Thus, $S_1 = S_2$ for all $x_0 \leq x \leq x_2$, i.e. x_1 is not really a knot. Similarly, x_{n-1} is not really a knot.
 \Rightarrow This condition is the default for Matlab's spline function.

3. Periodic spline: Assume that the data is periodic, i.e. $y_0 = y_n$ ($S_1(a) = S_n(b)$).

$$\begin{aligned} S_1'(a) &= S_n'(b) \\ S_1''(a) &= S_n''(b) \end{aligned}$$

4. Clamped spline: $S_1'(a) = v_0$, $S_n'(b) = v_n$, where $v_0, v_n \in \mathbb{R}$ are given.

Theorem 17.5.

For each of the 4 cases in Remark 17.4, there exists a unique cubic spline that interpolates the data.

Remark 17.6.

Cubic splines are the “smoothest” interpolants among all interpolating functions $f \in C^2[a, b]$.

$$\|f\|_2 = \left(\int_a^b |f(x)|^2 dx \right)^{1/2}$$

Theorem 17.7.

Let S be a cubic spline that interpolates the data. Let $f \in C^2[a, b]$ be any function that also interpolates the data: $f(x_j) = y_j$, $j = 0, 1, \dots, n$. Assume that

$$S''(x)[f'(x) - S'(x)]|_{x=a}^{x=b} = 0.$$

Then $\|S''\|_2 \leq \|f''\|_2$.

Remark 17.8. Comment on Theorem 17.7

Theorem 17.7 does not apply for the not-a-knot condition because the assumption does not hold, but it does apply for the other 3.

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18.1 Proof of Theorem 17.7

Proof.

$$\begin{aligned}
\int_a^b S''(x)[f''(x) - S''(x)] dx &\stackrel{\text{IBP}}{=} \underbrace{\int_a^b S''(f' - S')|_a^b}_{=0} - \int_a^b S'''(x)[f'(x) - S'(x)] dx \\
&= - \sum_{j=1}^n c_j \int_{x_{j-1}}^{x_j} f'(x) - S'(x) dx & S'''(x) = c_j = \text{constant on } [x_{j-1}, x_j] \\
&= - \sum_{j=1}^n c_j [f(x) - S(x)]|_{x_{j-1}}^{x_j} \\
&= - \sum_{j=1}^n c_j (y_j - y_{j-1} - y_{j-1} + y_{j+1}) \\
&= 0
\end{aligned}$$

Next,

$$\begin{aligned}
f''(x) &= S''(x) + [f''(x) - S''(x)] \\
[f''(x)]^2 &= [S''(x)]^2 + 2S''(x)[f''(x) - S''(x)] + [f''(x) - S''(x)]^2 \\
\|f''\|_2^2 &= \underbrace{\int_a^b [S''(x)]^2 dx}_{=\|S''\|_2^2} + 2 \underbrace{\int_a^b S''(x)[f''(x) - S''(x)] dx}_{=0} + \underbrace{\int_a^b [f''(x) - S''(x)]^2 dx}_{\geq 0} \\
\|f''\| &\geq \|S''\|_2
\end{aligned}$$

□

18.2 Construction of an Interpolating Cubic Spline

Constructing splines boils down to solving a tridiagonal linear system (\Rightarrow you never have to pivot). We will only cover the natural cubic spline; the other endpoint conditions are similar.

Recall:

$$a = x_0 < x_1 < x_2 < \cdots < x_n = b$$

$$S(x_j) = y_j, \quad j = 0, 1, \dots, n$$

$$S(x) = S_j(x) = \alpha_j + \beta_j(x - x_{j-1}) + \gamma_j(x - x_{j-1})^2 + \delta_j(x - x_{j-1})^3 \quad \text{on } [x_{j-1}, x_j], \quad j = 1, 2, \dots, n$$

Conditions:

1. $S_j(x_{j-1}) = y_{j-1}$, $S_j(x_j) = y_j$ for $j = 1, 2, \dots, n$
2. $S'_j(x_j) = S'_{j+1}(x_j)$, $j = 1, 2, \dots, n-1$
3. $S''_j(x_j) = S''_{j+1}(x_j)$, $j = 1, 2, \dots, n-1$
4. $S''_1(x_0) = 0$, $S''_n(x_n) = 0$ (natural spline condition)

From these conditions, we get:

1. $\alpha_j = y_{j-1}$ and

$$\begin{aligned} y_j - y_{j-1} &= \beta_j(x_j - x_{j-1}) + \gamma_j(x_j - x_{j-1})^2 + \delta_j(x_j - x_{j-1})^3 \\ \Delta_j := y_j - y_{j-1}, \quad \xi_j &:= x_j - x_{j-1} \end{aligned} \quad (\text{Notation}) \quad (18.1)$$

2. $\beta_j + 2\gamma_j\xi_j + 3\delta_j\xi_j^2 = \beta_{j+1}$, $j = 1, 2, \dots, n-1$

3. $2\gamma_j + 6\delta_j\xi_j = 2\gamma_{j+1}$, $j = 1, 2, \dots, n$

4. $2\gamma_1 = 0$, $2\gamma_n + 6\delta_n n \xi_n = 0 \Leftrightarrow$

$\gamma = 0$, $\gamma_{n+1} = 0$ where $\gamma_{n+1} = \gamma_n + 3\delta_n \xi_n$

3.

$$\delta_j = \frac{\gamma_{j+1} - \gamma_j}{3\xi_j}$$

1.

$$\beta_j = \frac{\Delta_j}{\xi_j} - \gamma_j \xi_j - \frac{\xi_j}{3}(\gamma_{j+1} - \gamma_j) = \frac{\Delta_j}{\xi_j} - \frac{\xi_j}{3}(\gamma_{j+1} + 2\gamma_j)$$

Insert the newest (3) and (1) into (2):

$$\xi_j \gamma_j + 2(\xi_j + \xi_{j+1})\gamma_{j+1} + \xi_{j+2}\gamma_{j+2} = 3 \left(\frac{\Delta_{j+1}}{\xi_{j+1}} - \frac{\Delta_j}{\xi_j} \right), \quad j = 1, 2, \dots, n-1$$

(where $\gamma_1 = 0$, $\gamma_{n+1} = 0$).

$$\left[\begin{array}{ccc|c} 2(\xi_1 + \xi_2) & \xi_2 & & \gamma_2 \\ \xi_2 & 2(\xi_2 + \xi_3) & \xi_3 & \gamma_3 \\ \ddots & \ddots & \ddots & \vdots \\ \ddots & \ddots & \ddots & \gamma_n \\ \ddots & \ddots & \ddots & \\ \xi_{n-2} & 2(\xi_{n-2} + \xi_{n-1}) & \xi_{n-1} & t_1 \\ \xi_{n-1} & 2(\xi_{n-1} + \xi_n) & 2(\xi_n + \xi_1) & t_2 \\ & & & \vdots \\ & & & t_{n-1} \end{array} \right] = \left[\begin{array}{c} t_1 \\ t_2 \\ \vdots \\ t_{n-1} \end{array} \right]$$

where

$$t_j := 3 \left(\frac{\Delta_{j+1}}{\xi_{j+1}} - \frac{\Delta_j}{\xi_j} \right), \quad j = 1, 2, \dots, n-1.$$

Remark 18.1.

The coefficient matrix of this system is tridiagonal, symmetric, and strictly diagonally dominant:

$$2(\xi_j + \xi_{j+1}) > \xi_j + \xi_{j+1} \quad \text{for all } j.$$

\Rightarrow There exists a unique solution $\gamma_2, \gamma_3, \dots, \gamma_n$.

Furthermore, no pivoting is needed. The system can be solved in $O(n)$ flops.

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19.1 B-splines

Given: $\Delta = \{x_0, x_1, \dots, x_n\}$, where $a = x_0 < x_1 < \dots < x_n = b$ and $n \geq 1$ (so that $a < b$). Recall that for the spline S of degree $k - 1$ (order k): $S \in C^{k-2}[a, b]$ and $S \in \prod_{k-1}$ on $[x_{j-1}, x_j]$, $j = 1, 2, \dots, n$.

Definition 19.1. $\mathcal{S}_{k,\Delta}$

$\mathcal{S}_{k,\Delta} =$ the set of all splines of order k , $k \geq 2$.

It is easy to see that $\mathcal{S}_{k,\Delta}$ is a real vector space:

$$S, \tilde{S} \in \mathcal{S}_{k,\Delta}, \alpha \in \mathbb{R} \Rightarrow S + \tilde{S} \in \mathcal{S}_{k,\Delta}, \alpha S \in \mathcal{S}_{k,\Delta}.$$

Also, $\prod_{k-1} \subset \mathcal{S}_{k,\Delta}$.

Proposition 19.2. *Dimension of $\mathcal{S}_{k,\Delta}$*

- On $[x_0, x_1]$: k coefficients
- On $[x_1, x_2]$: k coefficients and $k - 1$ conditions \Rightarrow i.e. only 1 coefficient
- Similarly for $[x_{j-1}, x_j]$, $j = 1, 2, 3, \dots, n$.

So $\dim(\mathcal{S}_{k,\Delta}) = k + n - 1$.

19.1.1 B-spline Basis (1)

We consider this kind of basis composed of monomials – that is, $p(x) = x^l$, $l = 0, 1, 2, \dots$ – and truncated powers of degree $k - 1$:

$$f_{x_j}(x) = (x - x_{j-1})_+^{k-1} := (\max\{x - x_j, 0\})^{k-1} = \begin{cases} (x - x_j)^{k-1} & x \geq x_j \\ 0 & x < x_j \end{cases}$$

$$\frac{d^l}{dx^l} (x - x_j)_+^{k-1} = (k-1)(k-2)\dots(k-1-l)(x - x_j)_+^{k-1-l}$$

So $(x - x_j)_+^{k-1} \in C^{k-2}[a, b]$, $(x - x_j)_+^{k-1} \in \mathcal{S}_{k,\Delta}$, and

$$\frac{d^{k-1}}{dx^{k-1}} (x - x_j)_+^{k-1} = \begin{cases} 0 & x < x_j \\ (k-1)! & x > x_j \end{cases}$$

Theorem 19.3.

The set

$$B = \{1, x, x^2, \dots, x^{k-1}, (x - x_1)_+^{k-1}, (x - x_2)_+^{k-1}, \dots, (x - x_{n-1})_+^{k-1}\}$$

is a basis for $\mathcal{S}_{k,\Delta}$.

Proof. We only need to show that the $k + n - 1$ functions in B are linearly independent. Let

$$S(x) = \sum_{l=0}^{k-1} \alpha_l x^l + \sum_{i=1}^{n-1} \beta_i (x - x_i)_+^{k-1} = 0$$

for all $a \leq x \leq b$. For each $j = 1, 2, \dots, n - 1$,

$$\begin{aligned} 0 &= S^{(k-1)}(x_j + 0) - S^{(k-1)}(x_j - 0) \\ &= \beta_j \left(\frac{d^{k-1}}{dx^{k-1}} (x - x_j)_+^{k-1} \Big|_{x=x_j+0} - \frac{d^{k-1}}{dx^{k-1}} (x - x_j)_+^{k-1} \Big|_{x=x_j-0} \right) \\ &= \beta_j (k-1)! . \end{aligned}$$

Thus, $\beta_j = 0$ for $j = 1, 2, \dots, n - 1$, and

$$S(x) = \sum_{l=0}^{k-1} \alpha_l x^l = 0$$

for all $a \leq x \leq b$. So $\alpha_l = 0$ for $l = 0, 1, \dots, k - 1$, and thus the functions in B are linearly independent. \square

However, the basis B is useless in practice:

- Ill-conditioned
- Global support. $\text{supp } x^l = [a, b]$, $\text{supp } (x - x_j)_+^{n-1} = [x_j, b]$. It is more efficient to have a basis function with local support only.

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20.1 Homework Comments

Problem 3:

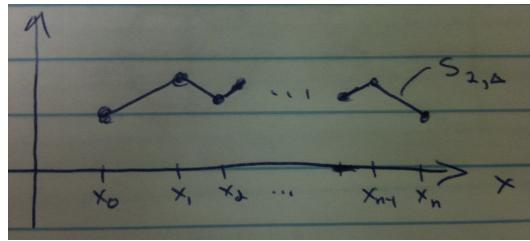
$$\frac{\|x - y\|}{\|x\|}$$

$Ax = b$, $b \neq 0$.

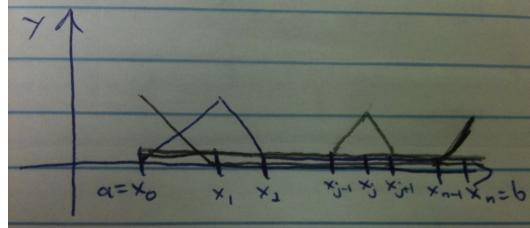
20.2 B-Splines

B-splines are basis functions for $\mathcal{S}_{k,\Delta}$ with desirable properties.

Example 20.1. $k = 2$ (Linear Splines)



B-splines for $\mathcal{S}_{2,\Delta}$ are the hat functions.



$$S(x) = \sum_{i=1}^{n+1} y_{i-1} N_{i2}(x)$$

$$S(x_{i-1}) = y_{i-1}, \quad i = 1, 2, \dots, n + 1$$

Note: the end points x_0 and x_n are special: for $k = 2$, we count x_0 and x_n “twice.”

For general $k \geq 2$:

$$\underbrace{x_0}_{\tau_1=\tau_2=\dots=\tau_k} \underset{k \text{ times}}{<} \underbrace{x_1}_{\tau_{k+1}} < \underbrace{x_2}_{\tau_{k+2}} < \dots < \underbrace{x_{n-1}}_{\tau_{k+n-1}} < \underbrace{x_n}_{\tau_{k+n}=\tau_{k+n+1}=\dots=\tau_{2k+n-1}}$$

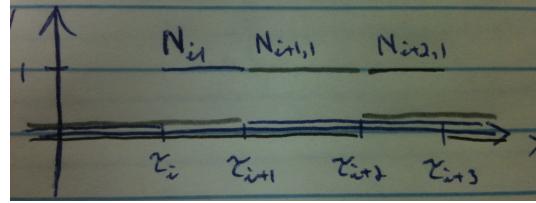
20.3 Construction of B-Splines of Order k (≥ 2)

$$\begin{aligned}\tau_1 &\leq \tau_2 \leq \dots \leq \tau_l, & \tau_1, \tau_2, \dots, \tau_l &\in \mathbb{R} \\ l &= 2k + n - 1 \\ l - k &= k + n - 1 \\ &= \dim \mathcal{S}_{k,\Delta}\end{aligned}\tag{20.1}$$

Definition 20.2.

The B-splines N_{ik} of order k (≥ 2) (associated with (20.1)), $i = 1, 2, \dots, l - k$, are defined recursively as follows:

$$N_{i1}(x) := \begin{cases} 1 & \tau_i \leq x < \tau_{i+1} \\ 0 & \text{otherwise} \end{cases}, \quad i = 1, 2, \dots, l - 1$$



For $j = 2, 3, \dots, k$:

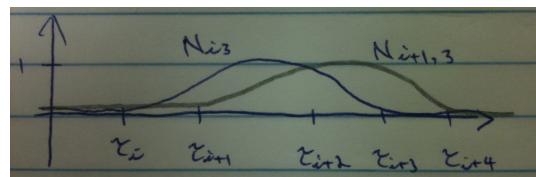
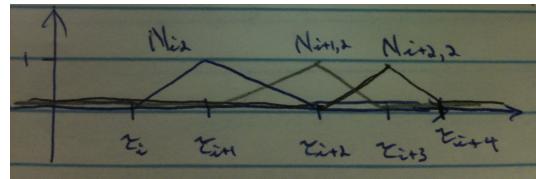
$$N_{ij}(x) := \frac{x - \tau_i}{\tau_{i+j-1} - \tau_i} N_{i,j-1}(x) + \frac{\tau_{i+j} - x}{\tau_{i+j} - \tau_{i+1}} N_{i+1,j-1}(x), \quad i = 1, 2, \dots, l - j \quad (20.2)$$

Conventions:

1. $\frac{0}{0} = 0$ if $\tau_{i+j} - \tau_i = 0$ or $\tau_{i+j} - \tau_{i+1} = 0$. For example, if $\tau_i = \tau_{i+1} = \dots = \tau_{i+k}$.

$$\begin{aligned} N_{i,j-1}(x) &\equiv 0 \\ N_{i+1,j-1}(x) &\equiv 0 \\ N_{ij}(x) &\equiv 0 \end{aligned}$$

2. $N_{ik}(\tau_l) := \lim_{x \rightarrow \tau_l^-} N_{ik}(x)$



Note:

$$\begin{aligned} N'_{i4}(\tau_i) &= N''_{i4}(\tau_i) = 0 \\ N'_{i4}(\tau_{i+4}) &= N''_{i4}(\tau_{i+4}) = 0 \end{aligned}$$

Properties

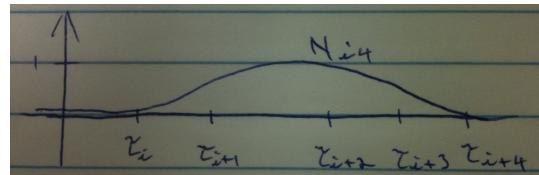
- (a) $\text{supp } N_{ik} \subset [\tau_i, \tau_{i+k}]$

(b) $N_{ik}(x) \geq 0$ for all $x \in [\tau_n, \tau_l]$

(c) $\sum_{i=1}^{l-k} N_{ik}(x) = 1$ for all $x \in [\tau_1, \tau_l]$, i.e. the N_{ik} 's form a partition of unity on $[\tau_1, \tau_l]$

(d) $N_{ik} \in C^{k-2}[\tau_1, \tau_l]$

(e) $N_{ik} \in \prod_{k-1}$ on $[\tau_j, \tau_{j+1}]$, $j = 1, 2, \dots, l - 1$

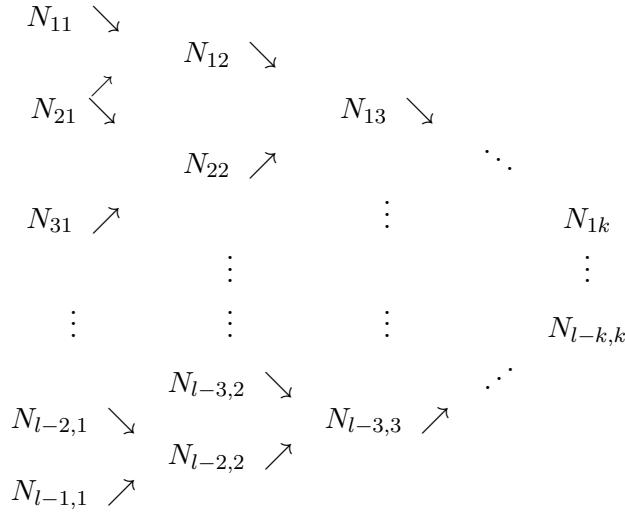


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21.1 B-Splines (Continued)

$$\tau_1 \leq \tau_2 \leq \cdots \leq \tau_l \quad (21.1)$$

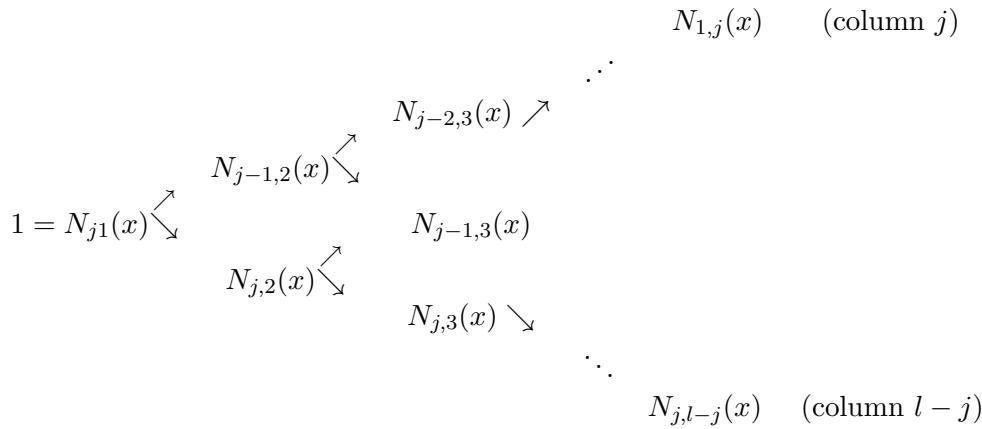
$$\text{For } j = 2, 3, \dots, k : \quad N_{ij}(x) = \cdots N_{i,\gamma-1}(x) + \cdots N_{i+1,j-1}(x), \quad i = 1, 2, \dots, l-j \quad (21.2)$$



21.1.1 Efficient Evaluation for $\tau_1 \leq x < \tau_l$

Determine j such that $\tau_j \leq x < \tau_{j+1}$ ($1 \leq j < l$, j is unique).

$$\begin{aligned}
N_{j1}(x) &= 1 \\
N_{ji}(x) &= 0 \quad \text{if } i \neq j
\end{aligned}$$



21.1.2 Back to $S_{k,\Delta}$

Choose the τ_j 's (21.1) as follows:

$$\begin{array}{ccccccc}
\underbrace{a}_{=\tau_1=\tau_2=\cdots=\tau_k} & = x_0 & < & \underbrace{x_1}_{=\tau_{k+1}} & < & \underbrace{x_2}_{=\tau_{k+2}} & < \cdots < \underbrace{x_{n-1}}_{=\tau_{k+n-1}} & < x_n = & \underbrace{b}_{=\tau_{k+n}=\tau_{k+n+1}=\cdots=\tau_{n+2k-1}} \\
& & & & & & & & & k \text{ times}
\end{array}$$

Set $l := n + 2k - 1$ ($\Rightarrow l - k = n + k - 1 = \dim \mathcal{S}_{k,\Delta}$).

Corresponding B-splines of order k :

$$N_{1k}, N_{2k}, \dots, N_{n+k-1,k}$$

Theorem 21.1.

The B-splines N_{ik} , $i = 1, 2, \dots, n + k - 1$ form a basis of $\mathcal{S}_{k,\Delta}$.

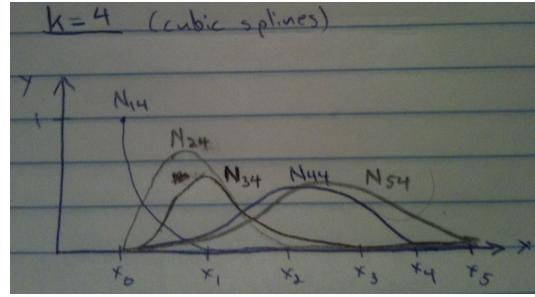


Figure 4: $\text{supp } N_{i4} \subseteq [\tau_i, \tau_{i+4}]$. Also, $\sum_i N_{i4}(x) = 1$.

Notes:

1. The representation

$$S(x) = \sum_{i=1}^{n+k-1} \gamma_i N_{ik}(x)$$

of any $S \in \mathcal{S}_{k,\Delta}$ is well-conditioned.

2. The coefficients γ_i , $i = 1, 2, \dots, n + k - 1$ are called the *de Boor points* of S .

3. Since $N_{ik}(x) \geq 0$ and $\sum_{i=1}^{n+k-1} N_{ik}(x) = 1$ for all $a \leq x \leq b$, the $S(x) = \sum_{i=1}^{n+k-1} \gamma_i N_{ik}(x)$ is a convex combination of the de Boor points $\gamma_1, \gamma_2, \dots, \gamma_{l-k}$ for any fixed $x \in [a, b]$.

Example 21.2. Cubic Spline with Not-A-Knot Condition

$k = 4, n \geq 3$.

$$\underbrace{x_0}_{=\tau_1=\tau_2=\tau_3=\tau_4} < x_1 < \underbrace{x_2}_{=\tau_5} < \cdots < \underbrace{x_{n-2}}_{=\tau_{n+1}} < x_{n-1} < \underbrace{x_n}_{=\tau_{n+2}=\tau_{n+3}=\tau_{n+4}=\tau_{n+5}}$$

$l = n + 5$. Corresponding B-splines $N_{14}, N_{24}, \dots, N_{n+1,4}$.

$$S(x) = \sum_{i=1}^{n+1} \gamma_i N_{i4}(x)$$

Interpolating Conditions:

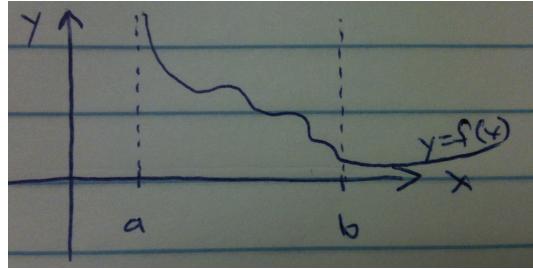
$$S(x_{j-1}) = y_{j-1}, \quad j = 1, \dots, n+1$$

$$M \begin{bmatrix} \gamma_1 \\ \gamma_2 \\ \vdots \\ \gamma_{n+1} \end{bmatrix} = \begin{bmatrix} y_0 \\ y_1 \\ \vdots \\ y_n \end{bmatrix}$$

where $M = [m_{ji}]_{j,i=1,2,\dots,n+1}$ with $m_{ji} = N_{i4}(x_{j-1})$.

22 11-14-11

22.1 Numerical Integration



Problem:

Evaluate

$$I = \int_a^b f(x) dx$$

where $f : (a, b) \rightarrow \mathbb{R}$ is integrable on $[a, b]$.

Quadrature Rules

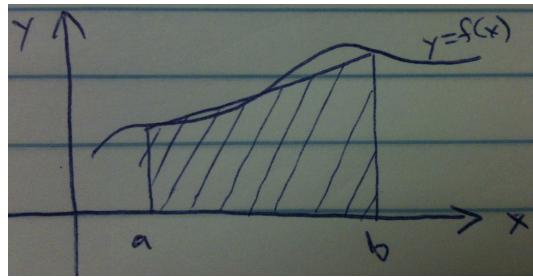
$$I \approx \sum_{i=1}^n w_i f(x_i)$$

where $a \leq x_1 < x_2 < \dots < x_n \leq b$ and $w_i \in \mathbb{R}$.

$$R_n := \sum_{i=1}^n w_i f(x_i) - I = (\text{remainder or error})$$

22.2 Examples of Quadrature Rules

1. Trapezoidal Rule



$$I \approx (b-a) \frac{f(a) + f(b)}{2} = \frac{b-a}{2} f(a) + \frac{b-a}{2} f(b)$$

$$n = 2, \quad x_1 = a, \quad x_2 = b, \quad w_1 = w_2 = \frac{b-a}{2}.$$

$$R_2 = (b-a)^3 \frac{1}{12} f^{(2)}(\xi) \quad \text{for some } \xi \in (a, b)$$

Thus, the trapezoid rule is exact for polynomials of degree ≤ 1 .

2. Simpson's Rule

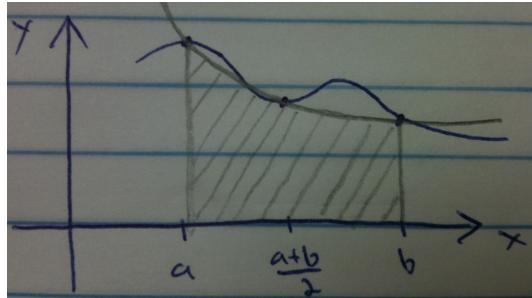


Figure 5: Interpolate the 3 points with a quadratic function.

$$I \approx \frac{b-a}{6} \left(f(a) + 4f\left(\frac{a+b}{2}\right) + f(b) \right)$$

$$n = 3, \quad x_1 = a, \quad x_2 = \frac{a+b}{2}, \quad x_3 = b, \quad w_1 = w_3 = \frac{b-a}{6}, \quad w_2 = \frac{2}{3}(b-a).$$

$$R_3 := (b-a)^5 \frac{1}{2880} f^{(4)}(\xi) \quad \text{for some } \xi \in (a, b)$$

Thus, Simpson's rule is exact for polynomials of degree ≤ 3 . Where does this extra degree of accuracy come from?

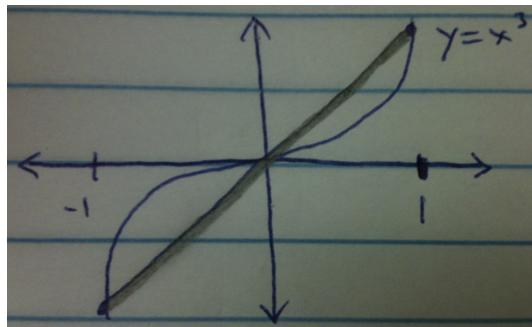
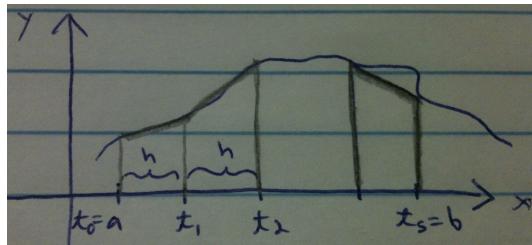


Figure 6: Simpson's rule for $y = x^3$ on $[-1, 1]$ is exact.

3. Compound trapezoidal rule. Let $s \geq 1$ be an integer. Set $t_i = a + ih$, $i = 0, 1, \dots, s$, where $h := \frac{b-a}{s}$.



$$\begin{aligned}
I &= \int_a^b f(x) dx = \sum_{i=0}^{s-1} \int_{t_i}^{t_{i+1}} f(x) dx \\
&\approx \sum_{i=0}^{s-1} \frac{h}{2} [f(t_i) + f(t_{i+1})] \\
&= \frac{h}{2} [f(a) + 2f(a+h) + 2f(a+2h) + \cdots + 2f(b-h) + f(b)] \\
&= h \left[\frac{1}{2} f(a) + f(a+h) + f(a+2h) + \cdots + f(b-h) + \frac{1}{2} f(b) \right] \\
&=: T_{s+1}
\end{aligned}$$

where $s+1$ is the number of points used. Error:

$$\begin{aligned}
R_{s+1} &= T_{s+1} - I = sh^3 \frac{1}{12} f^{(2)}(\xi) \quad \text{for some } \xi \in (a, b) \\
&= (b-a) \frac{1}{12} h^2 f^{(2)}(\xi)
\end{aligned}$$

T_{s+1} involves $s+1$ points ($h = \frac{b-a}{s}$).

T_{2s+1} involves $2s+1$ points ($\hat{h} = \frac{b-a}{2s} = \frac{h}{2}$).

$$\begin{aligned}
T_{2s+1} &= \underbrace{\hat{h}}_{\frac{h}{2}} \left[\frac{1}{2} f(a) + f(a+\hat{h}) + \underbrace{f(a+2\hat{h}) + \cdots + f(b-2\hat{h})}_{f(a+h)} + f(b-\hat{h}) + \frac{1}{2} f(b) \right] \\
&= \frac{1}{2} T_{s+1} + \hat{h} \left[f(a+\hat{h}) + f(a+3\hat{h}) + \cdots + f(b-3\hat{h}) + f(b-\hat{h}) \right]
\end{aligned}$$

Consequence: Once we have T_{s+1} , the approximation T_{s+1} can be obtained with s additional function evaluations! We can also estimate the error of T_{2s+1} :

$$\begin{aligned}
R_{s+1} &= T_{s+1} - I = (b-a) \frac{1}{12} h^2 f^{(2)}(\xi) \\
R_{2s+1} &= T_{2s+1} - I = (b-a) \frac{1}{12} \frac{h^2}{4} f^{(2)}(\hat{\xi})
\end{aligned}$$

Assume that $f^{(2)}(\xi) \approx f^{(2)}(\hat{\xi})$. Then

$$R_{2s+1} \approx \frac{1}{4} R_{s+1}.$$

But:

$$\begin{aligned}
|T_{2s+1} - T_{s+1}| &= |R_{2s+1} - R_{s+1}| \\
&= \frac{3}{4} |R_{s+1}|
\end{aligned}$$

So an error estimate for $T_{s+1} \approx |R_{s+1}| \approx \frac{4}{3} |T_{2s+1} - T_{s+1}|$. Approximate integral:

$$I \approx T_{2s+1}.$$

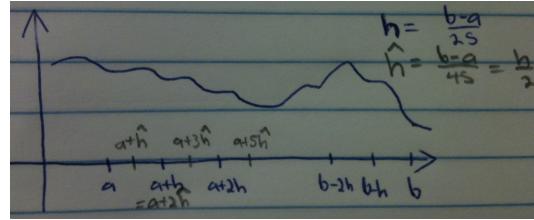
Conservative error estimate:

$$\frac{4}{3} |T_{2s+1} - T_{s+1}|.$$

23 11-16-11

23.1 Examples of Quadratures (Continued)

4. Compound Simpson's Rule. $s \geq 1$, $h = \frac{b-a}{2s}$



$$\begin{aligned}
I &= \int_a^b f(x) dx = \sum_{i=0}^{s-1} \int_{a+2ih}^{a+2(i+1)h} f(x) dx \\
&\approx \frac{h}{3} \left[\underbrace{f(a) + 4f(a+h) + f(a+2h)}_{+ f(b-2h) + 4f(b-h) + f(b)} + \underbrace{f(a+2h) + 4f(a+3h) + f(a+4h)}_{+ f(b-2h) + 4f(b-h) + f(b)} + \cdots + f(b-2h) \right] \\
&\approx \frac{h}{3} [f(a) + 4f(a+h) + 2f(a+2h) + 4f(a+3h) + 2f(a+4h) + \cdots + 2f(b-2h) + 4f(b-h) + f(b)] \\
&= S_{2s+1}
\end{aligned}$$

This uses $2s + 1$ points. Remainder:

$$R_{2s+1} = S_{2s+1} - I = (b-a) \frac{h^4}{180} f^{(4)}(\xi) \quad \text{for some } \xi \in (a, b)$$

If we double s :

$$R_{4s+1} \approx \frac{1}{16} R_{2s+1}$$

Approximate integral: $I \approx S_{4s+1}$

Conservative error estimate: $\frac{16}{15} |S_{4s+1} - S_{2s+1}|$

Efficient Implementation

$$\begin{aligned}
S_{2s+1} &= \frac{h}{3} \left[f(a) + 4 \sum_{j=1}^s f(a + (2j-1)h) + f(b) + 2 \sum_{k=1}^{s-1} \sum f(a + 2kh) \right], \quad h = \frac{b-a}{2s} \\
&= \frac{h}{3} \left[\underbrace{f(a) + 2 \sum_{i=1}^{2s-1} f(a + ih) + f(b)}_{=: A_{2s+1}} + \underbrace{2 \sum_{j=1}^s f(a + (2j-1)h)}_{=: B_s} \right] \\
S_{4s+1} &= \frac{h}{6} [A_{4s+1} + B_{2s}]
\end{aligned}$$

where $A_{4s+1} = f(a) + 2 \sum_{i=1}^{4s-1} f\left(a + i\frac{h}{2}\right) + f(b)$

and $B_{2s} = 2 \sum_{j=1}^{2s} f\left(a + (2j-1)\frac{h}{2}\right)$

We get A_{4s+1} for free because

$$\begin{aligned} A_{4s+1} &= f(a) + 2 \underbrace{\sum_{\substack{k=1 \\ (2k=i)}}^{2s-1} f(a + kh) + f(b)}_{=A_{2s+1}} + 2 \sum_{\substack{j=1 \\ (2j-1=i)}}^{2s} f\left(a + (2j-1)\frac{h}{2}\right) \\ &= A_{2s+1} + B_{2s} \end{aligned}$$

Consequence: To obtain S_{4s+1} from $S_{2s+1} = \frac{b-a}{6s}(A_{2s+1} + B_s)$, we only need to compute

$$B_{2s} = 2 \sum_{j=1}^{2s} f\left(a + (2j-1)\hat{h}\right), \quad \hat{h} = \frac{h}{2}$$

and set

$$\begin{aligned} A_{4s+1} &= A_{2s+1} + B_{2s} \\ S_{4s+1} &= \frac{b-a}{12s}(A_{4s+1} + B_{2s}) \end{aligned}$$

5. Gaussian Integration. So far all the formulas we've seen have been of the form

$$I = \int_a^b f(x) dx \approx \sum_{i=1}^n w_i f(x_i) \tag{23.1}$$

This expression has $2n$ degrees of freedom, so we hope to be able to integrate exactly polynomials of degree $\leq 2n - 1$. For Gaussian integration, we choose $a < x_1 < x_2 < \dots < x_n < b$ to be the zeros of the n th *orthogonal polynomial* $p_n(x) = x^n + \dots \in \prod_n$ defined by

$$\int_a^b p(x)p_n(x) dx = 0 \quad \text{for all } p \in \prod_n$$

($a = -1$, $b = 1$, $p_n(x) = \gamma_n \cdot (\text{nth Legendre polynomial})$.)

The w_i 's are chosen such that (23.1) is exact for all $f \in \prod_{n-1}$, i.e.

$$\int_a^b f(x) dx = \sum_{i=1}^n w_i f(x_i), \quad f \in \prod_n$$

(For example: $f(x) = x^j$, $j = 0, 1, \dots, n-1$.)

In this case, we get a linear system:

$$\begin{bmatrix} 1 & 1 & \cdots & 1 \\ x_1 & x_2 & \cdots & x_n \\ x_1^2 & x_2^2 & \cdots & x_n^2 \\ \vdots & \vdots & \ddots & \vdots \\ x_1^{n-1} & x_2^{n-1} & \cdots & x_n^{n-1} \end{bmatrix} \begin{bmatrix} w_1 \\ w_2 \\ \vdots \\ w_{n-1} \end{bmatrix} = \begin{bmatrix} b-a \\ \frac{1}{2}(b^2 - a^2) \\ \vdots \\ \frac{1}{n}(b^n - a^n) \end{bmatrix}$$

24 11-18-11

24.1 Gaussian Integration (Continued)

$$I = \int_a^b f(x) dx \approx \sum_{i=1}^n w_i f(x_i) \quad (24.1)$$

Choose the x_i as the zeros of $f'(x)$. $a < x_1 < x_2 < \dots < x_n < b$. The w_i 's are chosen so that

$$\int_a^b f(x) dx = \sum_{i=1}^n w_i f(x_i)$$

for all $f \in \prod_{n-1}$.

Theorem 24.1.

(24.1) is exact for all $p \in \prod_{2n-1}$.

Proof. Let $p \in \prod_{2n-1}$. $0 = \int_a^b p_n(x)q(x) dx$ for all $q \in \prod_{n-1}$. Polynomial division:

$$\begin{aligned} p(x) &= p_n(x)q(x) + r(x) \quad \text{where } q, r \in \prod_{n-1} \\ \int_a^b p(x) dx &= \underbrace{\int_a^b p_n(x)q(x) dx}_{=0 \text{ b/c}} + \int_a^b r(x) dx \\ &\quad p_n = \text{nth orthogonal polynomial} \\ &= \sum_{i=1}^n w_i r(x_i) \\ &= \sum_{i=1}^n w_i \underbrace{[p_n(x_i) q(x_i) + r(x_i)]}_{=p(x_i)} \\ &= \sum_{i=1}^n w_i p(x_i) \end{aligned}$$

□

24.1.1 Pros and Cons of Gaussian Integration

Pros

- Optimal Accuracy
- Can be used for functions with singularities at $x = a$ or $x = b$ (since $a < x_1 < x_2 < \dots < x_n < b$)

Cons

- When n is increased, old values $f(x_1), f(x_2), \dots, f(x_n)$ cannot be reused

24.1.2 Gauss-Kronrod Rules

Given: Gaussian rule with n points,

$$G_n = \sum_{i=1}^n w_i f(x_i).$$

Gauss-Kronrod:

$$K_{2n+1} = \sum_{i=1}^n a_i f(x_i) + \sum_{j=1}^{n+1} b_j f(y_j)$$

where the y_j 's are chosen as the zeros of the polynomial $q_{n+1}(x) = x^{n+1} + \dots \in \prod_{n+1}$ satisfying

$$\int_a^b p_n(x)p(x)q_{n+1}(x) dx = 0 \quad \text{for all } p \in \prod_n.$$

Thus, q_{n+1} is the $(n+1)$ st orthogonal polynomial with respect to the “inner product” $(p, q) := \int_a^b p(x)q(x)p_n(x) dx$.

The remaining parameters a_i , $i = 1, 2, \dots, n$ and b_j , $j = 1, 2, \dots, n+1$ are determined such that

$$\int_a^b p(x) dx = \sum_{i=1}^n a_i p(x_i) + \sum_{j=1}^{n+1} b_j p(y_j) \quad \text{for all } p \in \prod_{2n}.$$

Theorem 24.2.

K_{2n+1} is exact for all $p \in \prod_{3n+1}$.

Proof. For $p \in \prod_{3n+1}$:

$$\begin{aligned} p &= (p_n q_{n+1})t + r, \quad r \in \prod_{2n}, \quad t \in \prod_n \\ \int_a^b p(x) dx &= \underbrace{\int_a^b p_n(x)q_{n+1}(x)t(x) dx}_{=0} + \int_a^b r(x) dx \\ &= \sum_{i=1}^n a_i r(x_i) + \sum_{j=1}^{n+1} b_j r(x_i) \\ &= \sum_{i=1}^n a_i p(x_i) + \sum_{j=1}^{n+1} b_j p(y_j) \end{aligned}$$

□

24.1.3 Practical Use

Pair (G_n, K_{2n+1}) . Approximate the integral:

$$K_{2n+1} \approx \int_a^b f(x) dx.$$

Heuristic error estimate:

$$(200|G_n - K_{2n+1}|)^{3/2}.$$

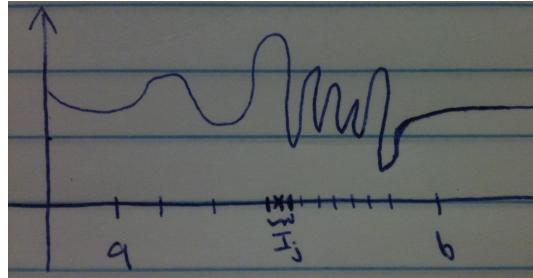
If we start with n reasonably large (e.g. $n = 10$), then K_{2n+1} will probably be accurate enough. However, if we need to double n then we have to start from scratch.

25 11-21-11

25.1 Corrections to the Homework

- 2(d) is trimmed down.
- use underscores for function names, not hyphens

25.2 Adaptive Quadrature



Idea: Use more function values $f(x_j)$ where f “varies more.”

Ingredients

- Quadrature rule Q with error estimate E .
- Tolerance $\epsilon > 0$.

Goal: Approximate the integral $Q_{[a,b]}$ such that

$$\left| Q_{[a,b]} - \int_a^b f(x) dx \right| \leq E_{[a,b]} < \epsilon.$$

25.2.1 Basic Adaptive Procedure

At every state:

$$[a, b] = I_1 \cup I_2 \cup \dots \cup I_l$$

where $I_j = [a_j, b_j]$, $h_j = b_j - a_j$. There is no overlap $\Rightarrow h_1 + h_2 + \dots + h_l = b - a$.

$$Q_j \approx \int_{I_j} f(x) dx, \quad \left| Q_j - \int_{I_j} f(x) dx \right| \leq E_j$$

Initialization

- $l = 1$, $I_1 = [a, b]$, $h_1 = b - a$. Apply your quadrature rule Q to $I_1 \rightarrow Q_1, E_1$.

General Step:

- If $E_j < \frac{h_j}{b-a}\epsilon$ for all $j = 1, 2, \dots, l$, stop. $Q_{[a,b]} = Q_1 + Q_2 + \dots + Q_l$, with error estimate

$$E_{[a,b]} = E_1 + E_2 + \dots + E_l < \frac{\epsilon}{b-a} \sum_{j=1}^l h_j = \epsilon$$

- Otherwise, for all $I_j = [a_j, b_j]$ with $E_j \geq \frac{h_j}{b-a}\epsilon$, set

$$I_j := \left[a_j, \frac{a_j + b_j}{2} \right], \quad I_{l+1} := \left[\frac{a_j + b_j}{2}, b_j \right]$$

$$h_j := h_{l+1} = \frac{b_j - a_j}{2}, \quad l := l + 1$$

and apply Q to I_j and $I_l \rightarrow Q_j, E_j$ and Q_l, E_l .

Note: For an actual algorithm, various safeguards are needed, e.g. $h_j \geq h_{\min} > 0$.

25.3 Eigenvalue Problems

Definition 25.1. Eigenvalue, Eigenvector

Let $A \in \mathbb{C}^{n \times n}$. A number $\lambda \in \mathbb{C}$ is called an *eigenvalue* of A if there exists a nonzero $x \in \mathbb{C}^n$, $x \neq \mathbf{0}$, such that

$$Ax = \lambda x.$$

Such an x is called an *eigenvector* of A .

Definition 25.2. Spectrum

The set

$$\Lambda(A) := \{\lambda \in \mathbb{C} \mid \lambda \text{ is an eigenvalue of } A\}$$

is called the *spectrum* of A .

$$\begin{aligned} Ax &= \lambda x, & x &\neq \mathbf{0} \\ \Rightarrow (\lambda I - A)x &= \mathbf{0}, & x &\neq \mathbf{0} \\ &\Rightarrow \text{the matrix } \lambda I - A \text{ is singular} \\ \Rightarrow \det(\lambda I - A) &= 0 \end{aligned}$$

Definition 25.3. Characteristic Polynomial

The polynomial

$$p(z) = p_A(z) := \det(zI - A) = z_n + \alpha_{n-1}z^{n-1} + \cdots + \alpha_1z + \alpha_0$$

is called the *characteristic polynomial* of A .

Notes:

1. The eigenvalues of A are the zeros of the characteristic polynomial, p_A . In particular, $A \in \mathbb{C}^{n \times n}$ has n eigenvalues, but they are not necessarily distinct.

2. For $n \geq 5$, any algorithm for computing eigenvalues of general $A \in \mathbb{C}^{n \times n}$ has to be iterative!

- Even if this was not true, working with polynomials is not great.

26 11-23-11

26.1 Eigenvalue Problems (Continued)

3. Any polynomial $p(z) = z^n + \alpha_{n-1}z^{n-1} + \alpha_{n-2}z^{n-2} + \cdots + \alpha_1z + \alpha_0$ is the characteristic polynomial of a matrix $A \in \mathbb{C}^{n \times n}$. For example,

$$A = \begin{bmatrix} 0 & 0 & \cdots & 0 & -\alpha_0 \\ 1 & \ddots & \ddots & \vdots & -\alpha_1 \\ 0 & \ddots & \ddots & 0 & \vdots \\ \vdots & \ddots & \ddots & 0 & -\alpha_{n-2} \\ 0 & \cdots & 0 & 1 & -\alpha_{n-1} \end{bmatrix}$$

has the characteristic polynomial

$$\det(zI - A) = \det \begin{bmatrix} z & 0 & \alpha_0 \\ -1 & z & \alpha_1 \\ & -1 & \ddots & \vdots \\ & & \ddots & z & \alpha_{n-2} \\ 0 & & & -1 & z + \alpha_{n-1} \end{bmatrix} = z^n + \alpha_{n-1}^{n-1} + \cdots + \alpha_1z + \alpha_0.$$

For some classes of matrices, the eigenvalues can be found trivially. For example, an upper triangular matrix:

$$A = \begin{bmatrix} a_{11} & * & \cdots & * \\ & a_{22} & \ddots & \vdots \\ & & \ddots & * \\ 0 & & & a_{nn} \end{bmatrix}, \quad \Lambda(A) = \{a_{11}, a_{22}, \dots, a_{nn}\}$$

26.2 Computation of Eigenvalues

26.2.1 Bad Ideas

1. Forming p_A and computing $\lambda \in \Lambda(A)$ as zeros of p_A .
2. Attempting to compute the *Jordan canonical form* of A :

$$X^{-1}AX = \begin{bmatrix} J_1 & & & 0 \\ & J_2 & & \\ & & \ddots & \\ 0 & & & J_l \end{bmatrix},$$

where each

$$J_i = \begin{bmatrix} \lambda_i & 1 & & 0 \\ & \lambda_i & 1 & \\ & & \ddots & \ddots & \\ 0 & & & \ddots & 1 \\ & & & & \lambda_i \end{bmatrix}, \quad i = 1, 2, \dots, l$$

are *Jordan blocks*.

`[J,X] = jordan(A)`

26.2.2 Better Ideas

1. Use the *Schur factorization* of A .
2. Use unitary *similarity transformations* to first transform A to a “simpler form.”

Definition 26.1. *Similarity Transformations*

Let $A \in \mathbb{C}^{n \times n}$. Let $U \in \mathbb{C}^{n \times n}$ be nonsingular. The map

$$A \mapsto U^{-1}AU$$

is called a *similarity transformation* of A .

Lemma 26.2.

$$\Lambda(A) = \Lambda(U^{-1}AU)$$

Proof.

$$\begin{aligned} Ax &= \lambda x, \quad x \neq \mathbf{0} \\ U^{-1}AU(\underbrace{U^{-1}x}_{\tilde{x}}) &= \lambda(\underbrace{U^{-1}x}_{\tilde{x}}), \quad \tilde{x} \neq 0 \\ U^{-1}AU\tilde{x} &= \lambda\tilde{x}, \quad \tilde{x} \neq 0 \end{aligned}$$

□

Use matrices U that are easy to invert and that are numerically well-behaved:

Definition 26.3. *Unitary*

$U \in \mathbb{C}^{n \times n}$ is said to be *unitary* if

$$U^H U = I$$

$$(U = [u_{jk}], \quad U^H = [\overline{u_{kj}}] = (\overline{U})^T.)$$

Notes:

1. U unitary $\Rightarrow U$ is nonsingular and $U^{-1} = U^H$.
2. U unitary $\Rightarrow U^H$ is unitary and $(U^H)^{-1} = (U^H)^H = U$.
3. U unitary $\Rightarrow \|Ux\|_2 = \|x\|_2$ for all $x \in \mathbb{C}^{n \times n}$. This is because $\|Ux\|_2^2 = (Ux)^H Ux = x^H U^H U x = x^H x = \|x\|_2^2$ and $\|U\|_2 = \max_{x \neq \mathbf{0}} \frac{\|Ux\|_2}{\|x\|_2} = 1$.

4. Unitary matrices U have the best possible Euclidean condition numbers:

$$\kappa_2(U) = \underbrace{\|U\|_2}_{=1} \underbrace{\|U^{-1}\|_2}_{=1}^{=U^H} = 1.$$

26.3 Unitary Similarity Transformations

$$A \mapsto U^{-1}AU = U^H AU =: T$$

where U is unitary.

Question: What is the “simplest” T we can achieve?

Theorem 26.4. Schur

For any $A \in \mathbb{C}^{n \times n}$, there exists a unitary $U \in \mathbb{C}^{n \times n}$ such that

$$U^H AU = T = \begin{bmatrix} t_{11} & * & \cdots & * \\ & t_{22} & \ddots & \vdots \\ & & \ddots & * \\ 0 & & & t_{nn} \end{bmatrix}$$

is upper-triangular. In particular, $\Lambda(A) = \Lambda(T) = \{t_{11}, t_{22}, \dots, t_{nn}\}$.

We will prove this by induction on n . For $n = 1$: $A = [a]$, $U = [1]$, $T = [a]$.

27 11-28-11

Office hours Wednesday 12:30-2:30.

27.1 Proof of Schur's Theorem

Recall:

$$U^H A U = T \text{ = upper-triangular}$$

Proof. By induction. For $n = 1$, it is trivial. Assume it is true for $1, \dots, n-1$. Let $\lambda \in \Lambda(A)$ with eigenvector x , $\|x\|_2 = 1$. Choose a unitary matrix

$$U_1 = [x \ * \ * \ \cdots \ *] \in \mathbb{C}^{n \times n}$$

with x as the first column. Then:

$$\begin{aligned} AU_1 &= [\lambda x \ * \ * \ \cdots \ *] \\ U_1^H A U_1 &= \begin{bmatrix} \lambda & * & \cdots & * \\ 0 & & & \\ \vdots & & \tilde{A} & \\ 0 & & & \end{bmatrix}, \quad \text{where } \tilde{A} \in \mathbb{C}^{n-1 \times n-1} \end{aligned}$$

Induction hypothesis:

$$\tilde{U}^H \tilde{A} \tilde{U} = \tilde{T} = \begin{bmatrix} * & * & \cdots & * \\ & * & \cdots & * \\ & & \ddots & \vdots \\ 0 & & & * \end{bmatrix}, \quad \tilde{U} \in \mathbb{C}^{n-1 \times n-1} \text{ is unitary}$$

Set

$$U := U_1 \begin{bmatrix} 1 & 0 & \cdots & 0 \\ 0 & & & \\ \vdots & & \tilde{U} & \\ 0 & & & \end{bmatrix} \in \mathbb{C}^{n \times n}$$

Then U is unitary and

$$\begin{aligned} U^H A U &= \begin{bmatrix} 1 & 0 & \cdots & 0 \\ 0 & & & \\ \vdots & & \tilde{U}^H & \\ 0 & & & \end{bmatrix} U_1^H A U_1 \begin{bmatrix} 1 & 0 & \cdots & 0 \\ 0 & & & \\ \vdots & & \tilde{U} & \\ 0 & & & \end{bmatrix} \\ &= \begin{bmatrix} \lambda & * & \cdots & * \\ 0 & & & \\ \vdots & & \tilde{U}^H \tilde{A} \tilde{U} & \\ 0 & & & \end{bmatrix} =: T = \text{upper-triangular} \end{aligned}$$

□

27.2 Two Simple Unitary Matrices: Householder Reflectors

Definition 27.1. Householder Reflectors

$Q \in \mathbb{C}^{n \times n}$ of the form $Q = I - 2vv^H$, where $v \in \mathbb{C}^n$ with $\|v\|_2 = 1$, is a *Householder reflector* matrix. Q is *Hermitian*:

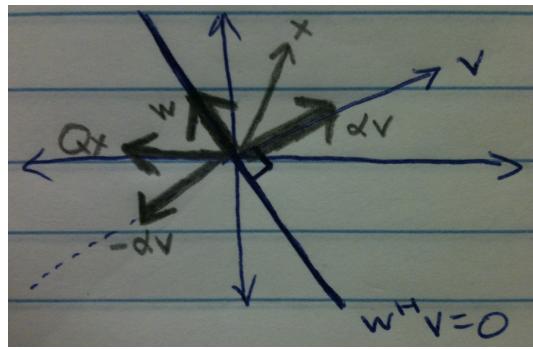
$$Q^H = I - 2(vv^H)^H = Q.$$

Q is unitary:

$$Q^H Q = Q^2 = I - 4vv^H + 4v \underbrace{v^H v}_{=\|v\|_2^2=1} v^H = I.$$

For any $x \in \mathbb{C}^n$:

Qx = reflection of x w.r.t. the $(n-1)$ -dimensional plane $\{w \in \mathbb{C}^n \mid w^H v = 0\}$



Proof.

$$\begin{aligned} x &= \alpha v + w, \quad \text{where } \alpha \in \mathbb{C}, \quad w^H v = 0 \\ Qx &= (I - 2vv^H)(\alpha v + w) \\ &= \alpha v + w - wv \underbrace{v^H \alpha v}_{= \alpha v^H H = \alpha} - 2v \underbrace{v^H w}_{= 0} \\ &= -\alpha v + w \end{aligned}$$

□

Given: $x \in \mathbb{C}^n$.

We can construct a Householder reflector $Q = I - 2vv^H$ such that

$$\begin{aligned} Qx &= \gamma e_1 \quad \text{for some } \gamma \in \mathbb{C} \text{ with } |\gamma| = \|x\|_2 \\ x &= \alpha v + w, \quad w^H v = 0 \\ Qx &= -\alpha v + w = \gamma e_1 \\ x - \gamma e_1 &= \alpha v + w - \gamma e_1 = 2\alpha v \end{aligned}$$

Set

$$\tilde{v} := x - \gamma e_1 = \begin{bmatrix} x_1 - \gamma \\ x_2 \\ \vdots \\ x_n \end{bmatrix} \quad \text{and} \quad v := \frac{\tilde{v}}{\|\tilde{v}\|_2}$$

Numerically best choice of γ :

$$\gamma := -(\operatorname{sgn} x_1) \|x\|_2$$

where

$$\operatorname{sgn} x_1 := \begin{cases} \frac{x_1}{\|x_1\|} & x_1 \neq 0 \\ 1 & x_1 = 0 \end{cases}.$$

Summary:

$$\begin{aligned} Q &= I - 2vv^H, \quad \text{where } v = \frac{\tilde{v}}{\|\tilde{v}\|_2} \\ \tilde{v} &= x + (\operatorname{sgn} x_1) \|x\|_2 e_1 \\ &= \begin{bmatrix} x_1 + (\operatorname{sgn} x_1) \|x\|_2 \\ x_2 \\ \vdots \\ x_n \end{bmatrix} \\ Qx &= \gamma e_1, \quad \text{where } \gamma = -(\operatorname{sgn} x_1) \|x\|_2 \end{aligned}$$

27.2.1 Reduction of A to Hessenberg Form

Definition 27.2. *Hessenberg Matrix*

http://en.wikipedia.org/wiki/Hessenberg_matrix

An *upper Hessenberg matrix* has zero entries below the first subdiagonal, and a *lower Hessenberg matrix* has zero entries above the first superdiagonal. For example:

$$\begin{bmatrix} 1 & 4 & 2 & 3 \\ 3 & 4 & 1 & 7 \\ 0 & 2 & 3 & 4 \\ 0 & 0 & 1 & 3 \end{bmatrix}$$

is upper Hessenberg.

$$\begin{bmatrix} 1 & 2 & 0 & 0 \\ 5 & 2 & 3 & 0 \\ 3 & 4 & 3 & 7 \\ 5 & 6 & 1 & 1 \end{bmatrix}$$

is lower Hessenberg.

Schur factorization:

$$U^H A U = T = \begin{bmatrix} * & \cdots & * \\ & \ddots & \vdots \\ 0 & & * \end{bmatrix}$$

Next “best” thing:

$$U^H A U = H = \begin{bmatrix} * & * & \cdots & \cdots & * \\ * & * & \ddots & & \vdots \\ & * & \ddots & \ddots & \vdots \\ & & \ddots & \ddots & * \\ 0 & & & * & * \end{bmatrix}$$

Can be achieved with $n - 2$ Householder reflections

$$U = U_1 U_2 \dots U_{n-2}, \quad \text{where } U_j = \begin{bmatrix} I_j & 0 \\ 0 & Q_{n-j} \end{bmatrix}, \quad Q_{n-j} \in \mathbb{C}^{n-j \times n-j}, \text{ Householder reflector}$$

Example 27.3. $n = 5$

$$A = \begin{bmatrix} a_{11} & * \\ * & * \end{bmatrix}$$

where $x \in \mathbb{C}^4$, $Q_1 \in \mathbb{C}^{4 \times 4}$, $Q_1x = \gamma_1 e_1$, $U_1 = \begin{bmatrix} 1 & 0 \\ 0 & Q_1 \end{bmatrix} = U_1^H$.

$$U_1^H A U_1 = U_1 A U_1 = (U_1 A) U_1 = \begin{bmatrix} a_{11} & * \\ \gamma_1 e_1 & * \end{bmatrix} U_1$$

$$= \begin{bmatrix} a_{11} & \# & \# & \# & \# \\ \gamma_1 & \# & \# & \# & \# \\ 0 & \# & \# & \# & \# \\ 0 & \# & \# & \# & \# \\ 0 & \# & \# & \# & \# \end{bmatrix}$$

$$U_3^H U_2^H U_1^H A U_1 U_2 U_3 = \begin{bmatrix} * & * & * & * & * \\ * & * & * & * & * \\ 0 & * & * & * & * \\ 0 & 0 & * & * & * \\ 0 & 0 & 0 & * & * \end{bmatrix}$$

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28.1 Reduction of A to Hessenberg Form (Continued)

$$U^H A U = H, \quad U = U_1 U_2 \dots U_{n-2}, \quad A \in \mathbb{C}^{n \times n}$$

Notes

1. $\Lambda(A) = \Lambda(H)$. Eigenvectors transform as follows:

$$\begin{aligned} Ax = \lambda x, \quad x \neq 0 &\Leftrightarrow \underbrace{U^H A U}_{H} U^H x = \lambda U^H x \\ y = U^H x &\Leftrightarrow Hy = \lambda y, \quad y \neq 0 \end{aligned}$$

Here $y = U^H x$ and $x = Uy = U_1 U_2 \dots U_{n-2} y$,

$$U_j = \begin{bmatrix} I_j & 0 \\ 0 & Q_{n-j} \end{bmatrix}, \quad Q_{n-j} = I - 2v_{n-j}v_{n-j}^H, \quad v_{n-j} \in \mathbb{C}^{n-j}$$

2. Reduction of $A \in \mathbb{C}^{n \times n}$ to upper-Hessenberg form requires $\sim \frac{10}{3}n^3$ flops.

Left to do: Compute the eigenvalues of H .

28.2 The QR Algorithm

Definition 28.1. QR Factorization

QR factorization of $M \in \mathbb{C}^{n \times n}$:

$$M = QR, \quad \text{where } Q \in \mathbb{C}^{n \times n} \text{ is unitary and } R \text{ is upper-triangular}$$

This is done via Gram-Schmidt.

Remark 28.2. QR Algorithm (for eigenvalue computations)

Input: upper-Hessenberg matrix $H \in \mathbb{C}^{n \times n}$

- Set $A^{(0)} := H$
- For $k = 1, 2, \dots$
 - Choose a suitable “shift” $\mu_k \in \mathbb{C}$
 - Compute a QR factorization:

$$A^{(k-1)} - \mu_k I = Q^{(k)} R^{(k)}$$

- Set $A^{(k)} := R^{(k)} Q^{(k)} + \mu_k I$
- end (k)

Lemma 28.3.

1. $A^{(k)} = (Q^{(k)})^H A^{(k-1)} Q^{(k)}, k = 1, 2, \dots$
2. $\Lambda(H) = \Lambda(A^{(k)}), k = 0, 1, 2, \dots$

Proof.

2. By (1), $A^{(k)}, A^{(k-1)}, A^{(k-2)}, \dots, A^{(1)}, A^{(0)} = H$ are all similar.

1.

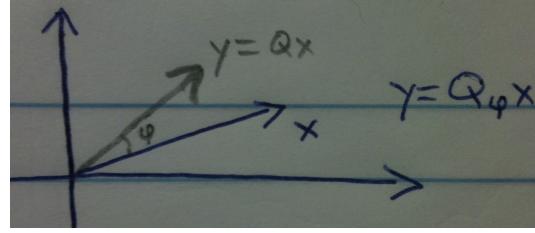
$$A^{(k)} = R^{(k)} Q^{(k)} + \mu_k I = (Q^{(k)})^H (\underbrace{Q^{(k)} R^{(k)} + \mu_k I}_{=A^{(k-1)}}) Q^{(k)}$$

□

Note: All the matrices $A^{(k)}$ are upper-Hessenberg!

28.3 Two Simple Unitary Matrices: Givens Rotations

28.3.1 2×2 Case



$$Q = \begin{bmatrix} c & s \\ -\bar{s} & \bar{c} \end{bmatrix} \in \mathbb{C}^{2 \times 2}, \quad \text{where } |c|^2 + |s|^2 = 1$$

Q is unitary:

$$Q^H Q = \begin{bmatrix} \bar{c} & -s \\ \bar{s} & c \end{bmatrix} \begin{bmatrix} c & s \\ -\bar{s} & \bar{c} \end{bmatrix} = \begin{bmatrix} |c|^2 + |s|^2 & 0 \\ 0 & |c|^2 + |s|^2 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}$$

Householder:

$$Q \begin{bmatrix} * \\ * \\ \vdots \\ * \end{bmatrix} = \begin{bmatrix} * \\ 0 \\ \vdots \\ 0 \end{bmatrix}$$

Givens:

$$Q \begin{bmatrix} * \\ * \end{bmatrix} = \begin{bmatrix} * \\ 0 \end{bmatrix}$$

Use of Q : to “zero out” the entry x_2 of any given $x = \begin{bmatrix} x_1 \\ x_2 \end{bmatrix} \in \mathbb{C}^2, x \neq 0$.

Indeed: Set

$$c = e^{i\alpha} \frac{\bar{x}_1}{\|x\|_2}, \quad s = e^{i\alpha} \frac{\bar{x}_2}{\|x\|_2}, \quad \text{where } \alpha \in \mathbb{C}$$

Then

$$|c|^2 + |s|^2 = \frac{|x_1|^2 + |x_2|^2}{\|x\|_2^2} = 1$$

$$Q \begin{bmatrix} x_1 \\ x_2 \end{bmatrix} = \frac{1}{\|x\|_2} \begin{bmatrix} e^{i\alpha} \bar{x}_1 & e^{i\alpha} \bar{x}_2 \\ -e^{-i\alpha} x_2 & e^{-i\alpha} x_1 \end{bmatrix} \begin{bmatrix} x_2 \\ x_2 \end{bmatrix} = \begin{bmatrix} e^{i\alpha} \|x\|_2 \\ 0 \end{bmatrix}$$

28.3.2 General Case

$$Q = \begin{bmatrix} 1 & & & & & 0 \\ & \ddots & & & & \\ & & 1 & & & \\ & & & c & & s \\ & & & & 1 & \\ & & & & & \ddots \\ & & & & & & 1 \\ & & & -\bar{s} & & \bar{c} & \\ & & & & & & 1 \\ 0 & & & & & & \ddots \\ & & & & & & & 1 \end{bmatrix}$$

Notes:

1. Left-multiplication with Q only changes the j th and k th rows.
2. Right-multiplication with Q only changes the j th and k th columns.
3. Q is a Givens rotation $\Rightarrow Q^H$ is a Givens rotation:

$$Q^H = \begin{bmatrix} \bar{c} & -s \\ \bar{s} & c \end{bmatrix}$$

28.3.3 Use in QR Algorithm

Example 28.4. $n = 5$

$$A^{(k-1)} - \mu_k I = \begin{bmatrix} * & * & * & * & * \\ * & * & * & * & * \\ 0 & * & * & * & * \\ 0 & 0 & * & * & * \\ 0 & 0 & 0 & * & * \end{bmatrix}$$

$$\text{Choose } Q_1 = \begin{bmatrix} c_1 & s_1 \\ -\bar{s}_1 & \bar{c}_1 \\ & & 1 \\ & & & 1 \\ & & & & 1 \end{bmatrix} \text{ such that } Q_1^H (A^{(k-1)} - \mu_k^T) = \begin{bmatrix} \# & \# & \# & \# & \# \\ 0 & \# & \# & \# & \# \\ 0 & * & * & * & * \\ 0 & 0 & * & * & * \\ 0 & 0 & 0 & * & * \end{bmatrix}$$

$$\text{Choose } Q_2 = \begin{bmatrix} 1 & & & & \\ & c_1 & s_1 & & \\ & -\bar{s}_1 & \bar{c}_1 & & \\ & & & 1 & \\ & & & & 1 \end{bmatrix} \text{ such that } Q_2^H Q_1^H (A^{(k-1)} - \mu_k^T) = \begin{bmatrix} \# & \# & \# & \# & \# \\ 0 & \# & \# & \# & \# \\ 0 & * & * & * & * \\ 0 & 0 & * & * & * \\ 0 & 0 & 0 & * & * \end{bmatrix}$$

⋮

$$Q_4^H Q_3^H Q_2^H Q_1^H (A^{(k-1)} - \mu_k^T) = \begin{bmatrix} * & * & * & * & * \\ 0 & * & * & * & * \\ 0 & 0 & * & * & * \\ 0 & 0 & 0 & \# & \# \\ 0 & 0 & 0 & 0 & \# \end{bmatrix} =: R^{(k)}$$

29 12-2-11

29.1 Comments on the Final

Friday 6-8 p.m.

- Office hours: Monday 12:30-2:30 p.m. and Wednesday 2-4 p.m.
- Open book & notes

29.2 QR Factorization (Continued)

For general $n \geq 2$

Convert A to a Hessenberg matrix and feed this to the QR algorithm: $A \rightarrow H = A^{(0)}$.

$$A^{(k-1)} - \mu_k I = Q^{(k)} R^{(k)} \quad \text{where } Q^{(k)} = \underbrace{Q_1 Q_2 \dots Q_{n-1}}_{n-1 \text{ Givens rotations}}$$

$$A^{(k)} = R^{(k)} Q^{(k)} + \mu_k I \quad \text{is upper-Hessenberg}$$

29.3 Convergence of the QR Algorithm

Let $\Lambda(A) = \{\lambda_1, \lambda_2, \dots, \lambda_n\}$ and assume that $|\lambda_1| > |\lambda_2| > \dots > |\lambda_n|$. (Special case: $\mu_k = 0$ for all k .) Then

$$A^{(k)} = \begin{bmatrix} * & * & \dots & \dots & * \\ a_{21}^{(k)} & * & \dots & \dots & * \\ & a_{32}^{(k)} & \ddots & & \vdots \\ & & \ddots & \ddots & \vdots \\ 0 & & & a_{n,n-1}^{(k)} & * \end{bmatrix} \xrightarrow{k \rightarrow \infty} R = \begin{bmatrix} r_{11} & * & \dots & \dots & * \\ 0 & r_{22} & \ddots & & \vdots \\ \vdots & \ddots & \ddots & \ddots & \vdots \\ \vdots & & \ddots & \ddots & * \\ 0 & \dots & \dots & 0 & r_{nn} \end{bmatrix}$$

where $\Lambda(R) = \{r_{11}, r_{22}, \dots, r_{nn}\} = \Lambda(A)$.

Speed of convergence: If $|\lambda_n| > 0$, then

$$a_{j,j-1}^{(k)} = O\left(\left|\frac{\lambda_0}{\lambda_{j-1}}\right|^k\right), \quad j = 2, 3, \dots, \quad k = 1, 2, \dots$$

\Rightarrow slow convergence if $|\lambda_j| \approx |\lambda_{j-1}|$. Shifts μ_k are used to speed up convergence. Suppose $\mu_k = \mu = \text{constant}$ and $|\lambda_1 - \mu| > |\lambda_2 - \mu| > \dots > |\lambda_n - \mu| > 0$. Speed of convergence:

$$a_{j,j-1}^{(k)} = O\left(\left|\frac{\lambda_j - \mu}{\lambda_{j-1} - \mu}\right|^k\right), \quad j = 2, 3, \dots, n, \quad k = 1, 2, \dots$$

$\Rightarrow a_{j,j-1}^{(k)} \xrightarrow{k \rightarrow \infty} 0$ fast if $\mu \approx \lambda_j$.

29.4 Strategy for Choosing μ_k

At the beginning of the k th iteration of the QR algorithm,

$$A^{(k-1)} = [a_{ij}^{(k-1)}] = \begin{bmatrix} * & \dots & \dots & \dots & * \\ a_{21}^{(k)} & * & & & \vdots \\ 0 & \ddots & \ddots & & 1 \\ \vdots & \ddots & \ddots & \ddots & \vdots \\ 0 & \dots & 0 & a_{n,n-1}^{(k)} & * \end{bmatrix}, \quad j = 2, 3, \dots, n, \quad \text{check if } a_{j,j-1}^{(k-1)} \approx 0.$$

$$a_{j-1,j-1}^{(k)} \\ a_{j,j-1}^{(k-1)} \quad a_{jj}^{(k)}$$

$$\left| a_{j,j-1}^{k-1} \right| \leq \epsilon \left(\left| a_{j-1,j-1}^{(k-1)} \right| + \left| a_{jj}^{(k-1)} \right| \right), \quad \text{where } \epsilon = O(e_{ps})$$

Set $a_{j,j-1}^{(k-1)} = 0$.

29.4.1 3 Cases

1. Case 1: $j = n$

$$A^{(k-1)} = \left[\begin{array}{cccccc|c} * & * & \cdots & \cdots & * & & * \\ * & \ddots & \ddots & & \vdots & & \vdots \\ & \ddots & \ddots & \ddots & \vdots & & \vdots \\ & & \ddots & \ddots & * & & \vdots \\ & & & * & * & & * \\ \hline 0 & \cdots & \cdots & 0 & \textcolor{red}{0} & a_{nn}^{(k-1)} & \end{array} \right] = \left[\begin{array}{c|c} \tilde{A}^{(k-1)} & * \\ \hline 0 & a_{nn}^{(k-1)} \end{array} \right], \quad \tilde{A}^{(k-1)} \in \mathbb{C}^{n-1 \times n-1}$$

$$\Rightarrow a_{nn}^{(k)} \in \Lambda(A^{(k-1)}) = \Lambda(A).$$

Accept $a_{nn}^{(k)}$ as an eigenvalue of A and continue the QR algorithm on $\tilde{A}^{(k-1)}$. Typical choice: $\mu_k = a_{n-1,n-1}^{(k-1)}$.

2. Case 2: $j = 2$

$$A^{(k-1)} = \left[\begin{array}{c|cccc} a_{11}^{(k-1)} & * & \cdots & \cdots & * \\ \hline \textcolor{red}{0} & & & & \\ 0 & & & & \\ \vdots & & \tilde{A}^{(k-1)} & & \\ 0 & & & & \end{array} \right]$$

$a_{11}^{(k-1)}$ is an eigenvalue of A .

3. Case 3: $2 < j < n$

$$A^{(k-1)} = \left[\begin{array}{c|ccccc|ccccc} * & * & \cdots & \cdots & * & * & \cdots & \cdots & \cdots & * \\ * & \ddots & \ddots & & \vdots & \vdots & & & & \vdots \\ & \ddots & \ddots & \ddots & \vdots & \vdots & & & & \vdots \\ & & \ddots & \ddots & * & \vdots & & & & \vdots \\ & & & * & * & * & \cdots & \cdots & \cdots & * \\ \hline & & & & \textcolor{red}{0} & * & * & \cdots & \cdots & * \\ & & & & & * & \ddots & \ddots & & \vdots \\ & & & & & & \ddots & \ddots & \ddots & \vdots \\ & & & & & & & \ddots & \ddots & * \\ 0 & & & & & & & & 0 & * & * \end{array} \right] = \left[\begin{array}{c|c} A_1^{(k-1)} & * \\ \hline 0 & A_2^{(k-1)} \end{array} \right]$$

$\Lambda(A) = \Lambda(A^{(k-1)}) = \Lambda(A_1^{(k-1)}) \cup \Lambda(A_2^{(k-1)})$. Continue QR on $A_1^{(k-1)} \in \mathbb{C}^{j-1 \times j-1}$ and $A_2^{(k-1)} \in \mathbb{C}^{n-j+1 \times n-j+1}$.

Typical flop count for such a practical QR algorithm:

- $\sim 10n^3$ if only the eigenvalues are computed
- $\sim 27n^3$ if eigenvalues and eigenvectors are computed

A Algorithms

Algorithm	Operation Count	Page
Triangular Solve	n^2	
LU Factorization without pivoting	$\frac{2n^3}{3}$	8
LU Factorization with partial pivoting		14
Newton's Method ($n = 1$)		17
Newton's Method		19
QR Algorithm	$10n^3$ ($27n^3$ with eigenvectors)	85
$A \in \mathbb{C}^{n \times n} \rightarrow$ upper-Hessenberg	$\frac{10n^3}{3}$	

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