## 3.7 Gaussian Elimination, LU-Factorization, and Cholesky Factorization

Let A be an  $n \times n$  matrix, let  $b \in \mathbb{R}^n$  be an n-dimensional vector and assume that A is invertible. Our goal is to solve the system Ax = b. Since A is assumed to be invertible, we know that this system has a unique solution,  $x = A^{-1}b$ . Experience shows that two counter-intuitive facts are revealed:

- (1) One should avoid computing the inverse,  $A^{-1}$ , of A explicitly. This is because this would amount to solving the n linear systems,  $Au^{(j)} = e_j$ , for j = 1, ..., n, where  $e_j = (0, ..., 1, ..., 0)$  is the jth canonical basis vector of  $\mathbb{R}^n$  (with a 1 is the jth slot). By doing so, we would replace the resolution of a single system by the resolution of n systems, and we would still have to multiply  $A^{-1}$  by b.
- (2) One does not solve (large) linear systems by computing determinants (using Cramer's formulae). This is because this method requires a number of additions (resp. multiplications) proportional to (n + 1)! (resp. (n + 2)!).

The key idea on which most direct methods (as opposed to iterative methods, that look for an approximation of the solution) are based is that if A is an upper-triangular matrix, which means that  $a_{ij} = 0$  for  $1 \le j < i \le n$  (resp. lower-triangular, which means that  $a_{ij} = 0$  for  $1 \le i < j \le n$ ), then computing the solution, x, is trivial. Indeed, say A is an upper-triangular matrix

$$A = \begin{pmatrix} a_{11} & a_{12} & \cdots & a_{1n-2} & a_{1n-1} & a_{1n} \\ 0 & a_{22} & \cdots & a_{2n-2} & a_{2n-1} & a_{2n} \\ 0 & 0 & \ddots & \vdots & \vdots & \vdots \\ & & \ddots & \vdots & \vdots \\ 0 & 0 & \cdots & 0 & a_{n-1n-1} & a_{n-1n} \\ 0 & 0 & \cdots & 0 & 0 & a_{nn} \end{pmatrix}.$$

Then,  $\det(A) = a_{11}a_{22}\cdots a_{nn} \neq 0$ , and we can solve the system Ax = b from bottom-up by back-substitution, i.e., first we compute  $x_n$  from the last equation, next plug this value of  $x_n$  into the next to the last equation and compute  $x_{n-1}$  from it, etc. This yields

$$x_{n} = a_{n n}^{-1} b_{n}$$

$$x_{n-1} = a_{n-1 n-1}^{-1} (b_{n-1} - a_{n-1 n} x_{n})$$

$$\vdots$$

$$x_{1} = a_{1 1}^{-1} (b_{1} - a_{1 2} x_{2} - \dots - a_{1 n} x_{n}).$$

If A was lower-triangular, we would solve the system from top-down by forward-substitution.

Thus, what we need is a method for transforming a matrix to an equivalent one in upper-triangular form. This can be done by *elimination*. Let us illustrate this method on the following example:

We can eliminate the variable x from the second and the third equation as follows: Subtract twice the first equation from the second and add the first equation to the third. We get the new system

$$2x + y + z = 5 
- 8y - 2z = -12 
8y + 3z = 14.$$

This time, we can eliminate the variable y from the third equation by adding the second equation to the third:

$$2x + y + z = 5$$
  
 $-8y - 2z = -12$   
 $z = 2$ 

This last system is upper-triangular. Using back-substitution, we find the solution: z = 2, y = 1, x = 1.

Observe that we have performed only row operations. The general method is to iteratively eliminate variables using simple row operations (namely, adding or subtracting a multiple of a row to another row of the matrix) while simultaneously applying these operations to the vector b, to obtain a system, MAx = Mb, where MA is upper-triangular. Such a method is called Gaussian elimination. However, one extra twist is needed for the method to work in all cases: It may be necessary to permute rows, as illustrated by the following example:

$$x + y + z = 1$$
  
 $x + y + 3z = 1$   
 $2x + 5y + 8z = 1$ .

In order to eliminate x from the second and third row, we subtract the first row from the second and we subtract twice the first row from the third:

$$x + y + z = 1$$
  
 $2z = 0$   
 $3y + 6z = -1$ .

Now, the trouble is that y does not occur in the second row; so, we can't eliminate y from the third row by adding or subtracting a multiple of the second row to it. The remedy is simple: Permute the second and the third row! We get the system:

which is already in triangular form. Another example where some permutations are needed is:

$$z = 1$$

$$-2x + 7y + 2z = 1$$

$$4x - 6y = -1$$

First, we permute the first and the second row, obtaining

and then, we add twice the first row to the third, obtaining:

$$\begin{array}{rclcrcr}
-2x & + & 7y & + & 2z & = & 1 \\
& & z & = & 1 \\
8y & + & 4z & = & 1.
\end{array}$$

Again, we permute the second and the third row, getting

$$\begin{array}{rclcrcr}
-2x & + & 7y & + & 2z & = & 1 \\
& & 8y & + & 4z & = & 1 \\
& & z & = & 1,
\end{array}$$

an upper-triangular system. Of course, in this example, z is already solved and we could have eliminated it first, but for the general method, we need to proceed in a systematic fashion.

We now describe the method of Gaussian Elimination applied to a linear system, Ax = b, where A is assumed to be invertible. We use the variable k to keep track of the stages of elimination. Initially, k = 1.

- (1) The first step is to pick some nonzero entry,  $a_{i1}$ , in the first column of A. Such an entry must exist, since A is invertible (otherwise, we would have  $\det(A) = 0$ ). The actual choice of such an element has some impact on the numerical stability of the method, but this will be examined later. For the time being, we assume that some arbitrary choice is made. This chosen element is called the *pivot* of the elimination step and is denoted  $\pi_1$  (so, in this first step,  $\pi_1 = a_{i1}$ ).
- (2) Next, we permute the row (i) corresponding to the pivot with the first row. Such a step is called *pivoting*. So, after this permutation, the first element of the first row is nonzero.
- (3) We now eliminate the variable  $x_2$  from all rows except the first by adding suitable multiples of the first row to these rows. More precisely we add  $-a_{i1}/\pi_1$  times the first row to the *i*th row, for i = 2, ..., n. At the end of this step, all entries in the first column are zero except the first.

(4) Increment k by 1. If k = n, stop. Otherwise, k < n, and then iteratively repeat steps (1), (2), (3) on the  $(n - k + 1) \times (n - k + 1)$  subsystem obtained by deleting the first k - 1 rows and k - 1 columns from the current system.

If we let  $A_1 = A$  and  $A_k = (a_{ij}^k)$  be the matrix obtained after k-1 elimination steps  $(2 \le k \le n)$ , then the kth elimination step is as follows: The matrix,  $A_k$ , is of the form

$$A_{k} = \begin{pmatrix} a_{11}^{k} & a_{12}^{k} & \cdots & \cdots & a_{1n}^{k} \\ & a_{22}^{k} & \cdots & \cdots & a_{2n}^{k} \\ & & \ddots & \vdots & & \vdots \\ & & a_{kk}^{k} & \cdots & a_{kn}^{k} \\ & & \vdots & & \vdots \\ & & a_{nk}^{k} & \cdots & a_{nn}^{k} \end{pmatrix}.$$

Now, we will prove later that  $\det(A_k) = \pm \det(A)$ . Since A is invertible, some entry  $a_{ik}^k$  with  $k \leq i \leq n$  is nonzero; so, one of these entries can be chosen as pivot, and we permute the kth row with the ith row, obtaining the matrix  $\alpha^k = (\alpha_{jl}^k)$ . The new pivot is  $\pi_k = \alpha_{kk}^k$ , and we zero the entries  $i = k + 1, \ldots, n$  in column k by adding  $-\alpha_{ik}^k/\pi_k$  times row k to row i. At the end of this step, we have  $A_{k+1}$ . Observe that the first k-1 rows of  $A_k$  are identical to the first k-1 rows of  $A_{k+1}$ .

It is easy to figure out what kind of matrices perform the elementary row operations used during Gaussian elimination. The permutation of the kth row with the ith row is achieved by multiplying A on the left by the transposition matrix, P(i,k), with

$$P(i,k)_{j\,l} = \begin{cases} 1 & \text{if } j=l, \text{ where } j \neq i \text{ and } l \neq k \\ 0 & \text{if } j=l=i \text{ or } j=l=k \\ 1 & \text{if } j=i \text{ and } l=k \text{ or } j=k \text{ and } l=i, \end{cases}$$

i.e.,

Observe that det(P(i,k)) = -1. Therefore, during the permutation step (2), if row k and row i need to be permuted, the matrix A is multiplied on the left by the matrix  $P_k$  such that  $P_k = P(i,k)$ , else we set  $P_k = I$ .

Adding  $\beta$  times row j to row i is achieved by multiplying A on the left by the elementary matrix,  $E_{i,j;\beta} = I + \beta e_{i,j}$ , where

$$(e_{i,j})_{kl} = \begin{cases} 1 & \text{if } k = i \text{ and } l = j \\ 0 & \text{if } k \neq i \text{ or } l \neq j, \end{cases}$$

i.e.,

Observe that the inverse of  $E_{i,j;\beta} = I + \beta e_{ij}$  is  $E_{i,j;-\beta} = I - \beta e_{i,j}$  and that  $\det(E_{i,j;\beta}) = 1$ . Therefore, during step 3 (the elimination step), the matrix A is multiplied on the left by a product,  $E_k$ , of matrices of the form  $E_{i,k;\beta_{i,k}}$ . Consequently, we see that

$$A_{k+1} = E_k P_k A_k$$
.

The fact that  $\det(P(i,k)) = -1$  and that  $\det(E_{i,j;\beta}) = 1$  implies immediately the fact claimed above: We always have  $\det(A_k) = \pm \det(A)$ . Furthermore, since

$$A_{k+1} = E_k P_k A_k$$

and since Gaussian elimination stops for k = n, the matrix

$$A_n = E_{n-1} P_{n-1} \cdots E_2 P_2 E_1 P_1 A$$

is upper-triangular. Also note that if we let  $M = E_{n-1}P_{n-1}\cdots E_2P_2E_1P_1$ , then  $\det(M) = \pm 1$ , and

$$\det(A) = \pm \det(A_n).$$

We can summarize all this in the following theorem:

**Theorem 3.14** (Gaussian Elimination) Let A be an  $n \times n$  matrix (invertible or not). Then there is some invertible matrix, M, so that U = MA is upper-triangular. The pivots are all nonzero iff A is invertible.

*Proof*. We already proved the theorem when A is invertible, as well as the last assertion. Now, A is singular iff some pivot is zero, say at stage k of the elimination. If so, we must have  $a_{ik}^k = 0$ , for  $i = k, \ldots, n$ ; but in this case,  $A_{k+1} = A_k$  and we may pick  $P_k = E_k = I$ .  $\square$ 

**Remark:** Obviously, the matrix M can be computed as

$$M = E_{n-1}P_{n-1}\cdots E_2P_2E_1P_1$$

but this expression is of no use. Indeed, what we need is  $M^{-1}$ ; when no permutations are needed, it turns out that  $M^{-1}$  can be obtained immediately from the matrices  $E_k$ 's, in fact, from their inverses, and no multiplications are necessary.

Remark: Instead of looking for an invertible matrix, M, so that MA is upper-triangular, we can look for an invertible matrix, M, so that MA is a diagonal matrix. Only a simple change to Gaussian elimination is needed. At every stage, k, after the pivot has been found and pivoting been performed if necessary, in addition to adding suitable multiples of the kth row to the rows below row k in order to zero the entries in column k for  $i = k + 1, \ldots, n$ , also add suitable multiples of the kth row to the rows above row k in order to zero the entries in column k for  $i = 1, \ldots, k - 1$ . Such steps are also achieved by multiplying on the left by elementary matrices  $E_{i,k;\beta_{i,k}}$ , except that i < k, so that these matrices are not lower-diagonal matrices. Nevertheless, at the end of the process, we find that  $A_n = MA$ , is a diagonal matrix. This method is called the Gauss-Jordan factorization. Because it is more expansive than Gaussian elimination, this method is not used much in practice. However, Gauss-Jordan factorization can be used to compute the inverse of a matrix, A. Indeed, we find the jth column of  $A^{-1}$  by solving the system  $Ax^{(j)} = e_j$  (where  $e_j$  is the jth canonical basis vector of  $\mathbb{R}^n$ ). By applying Gauss-Jordan, we are led to a system of the form  $D_j x^{(j)} = M_j e_j$ , where  $D_j$  is a diagonal matrix, and we can immediately compute  $x^{(j)}$ .

It remains to discuss the choice of the pivot, and also conditions that guarantee that no permutations are needed during the Gaussian elimination process.

We begin by stating a necessary and sufficient condition for an invertible matrix to have an LU-factorization (i.e., Gaussian elimination does not require pivoting). We say that an invertible matrix, A, has an LU-factorization if it can be written as A = LU, where U is upper-triangular invertible and L is lower-triangular, with  $L_{ii} = 1$  for  $i = 1, \ldots, n$ . A lower-triangular matrix with diagonal entries equal to 1 is called a unit lower-triangular matrix. Given an  $n \times n$  matrix,  $A = (a_{ij})$ , for any k, with  $1 \le k \le n$ , let A[1..k, 1..k] denote the submatrix of A whose entries are  $a_{ij}$ , where  $1 \le i, j \le k$ .

**Proposition 3.15** Let A be an invertible  $n \times n$ -matrix. Then, A, has an LU-factorization, A = LU, iff every matrix A[1..k, 1..k] is invertible for k = 1, ..., n.

*Proof.* First, assume that A = LU is an LU-factorization of A. We can write

$$A = \begin{pmatrix} A[1..k, 1..k] & A_2 \\ A_3 & A_4 \end{pmatrix} = \begin{pmatrix} L_1 & 0 \\ P & L_4 \end{pmatrix} \begin{pmatrix} U_1 & Q \\ 0 & U_4 \end{pmatrix} = \begin{pmatrix} L_1U_1 & L_1Q \\ PU_1 & PQ + L_4U_4 \end{pmatrix},$$

where  $L_1, L_4$  are unit lower-triangular and  $U_1, U_4$  are upper-triangular. Thus,

$$A[1..k, 1..k] = L_1U_1,$$

and since U is invertible,  $U_1$  is also invertible (the determinant of U is the product of the diagonal entries in U, which is the product of the diagonal entries in  $U_1$  and  $U_4$ ). As  $L_1$  is invertible (since its diagonal entries are equal to 1), we see that A[1..k, 1..k] is invertible for k = 1, ..., n.

Conversely, assume that A[1..k, 1..k] is invertible, for k = 1, ..., n. We just need to show that Gaussian elimination does not need pivoting. We prove by induction on k that the kth step does not need pivoting. This holds for k = 1, since  $A[1..1, 1..1] = (a_{11})$ , so,  $a_{11} \neq 0$ . Assume that no pivoting was necessary for the first k steps  $(1 \leq k \leq n - 1)$ . In this case, we have

$$E_{k-1}\cdots E_2E_1A=A_k,$$

where  $L = E_{k-1} \cdots E_2 E_1$  is a unit lower-triangular matrix and  $A_k[1..k, 1..k]$  is upper-triangular, so that  $LA = A_k$  can be written as

$$\begin{pmatrix} L_1 & 0 \\ P & L_4 \end{pmatrix} \begin{pmatrix} A[1..k, 1..k] & A_2 \\ A_3 & A_4 \end{pmatrix} = \begin{pmatrix} U_1 & B_2 \\ 0 & B_4 \end{pmatrix},$$

where  $L_1$  is unit lower-triangular and  $U_1$  is upper-triangular. But then,

$$L_1A[1..k, 1..k]) = U_1,$$

where  $L_1$  is invertible and  $U_1$  is also invertible since its diagonal elements are the first k pivots, by hypothesis. Therefore, A[1..k, 1..k] is also invertible.  $\square$ 

**Corollary 3.16** (LU-Factorization) Let A be an invertible  $n \times n$ -matrix. If every matrix A[1..k, 1..k] is invertible for k = 1, ..., n, then Gaussian elimination requires no pivoting and yields an LU-factorization, A = LU.

*Proof*. We proved in Proposition 3.15 that in this case Gaussian elimination requires no pivoting. Then, since every elementary matrix  $E_{i,k;\beta}$  is lower-triangular (since we always arrange that the pivot,  $\pi_k$ , occurs above the rows that it operates on), since  $E_{i,k;\beta}^{-1} = E_{i,k;-\beta}$  and the  $E'_k s$  are products of  $E_{i,k;\beta_{i,k}}$ 's, from

$$E_{n-1}\cdots E_2E_1A=U,$$

where U is an upper-triangular matrix, we get

$$A = LU$$
.

where  $L = E_1^{-1} E_2^{-1} \cdots E_{n-1}^{-1}$  is a lower-triangular matrix. Furthermore, as the diagonal entries of each  $E_{i,k;\beta}$  are 1, the diagonal entries of each  $E_k$  are also 1.  $\square$ 

The reader should verify that the example below is indeed an LU-factorization.

$$\begin{pmatrix} 2 & 1 & 1 & 0 \\ 4 & 3 & 3 & 1 \\ 8 & 7 & 9 & 5 \\ 6 & 7 & 9 & 8 \end{pmatrix} = \begin{pmatrix} 1 & 0 & 0 & 0 \\ 2 & 1 & 0 & 0 \\ 4 & 3 & 1 & 0 \\ 3 & 4 & 1 & 1 \end{pmatrix} \begin{pmatrix} 2 & 1 & 1 & 0 \\ 0 & 1 & 1 & 1 \\ 0 & 0 & 2 & 2 \\ 0 & 0 & 0 & 2 \end{pmatrix}.$$

One of the main reasons why the existence of an LU-factorization for a matrix, A, is interesting is that if we need to solve several linear systems, Ax = b, corresponding to the same matrix, A, we can do this cheaply by solving the two triangular systems

$$Lw = b$$
, and  $Ux = w$ .

As we will see a bit later, symmetric positive definite matrices satisfy the condition of Proposition 3.15. Therefore, linear systems involving symmetric positive definite matrices can be solved by Gaussian elimination without pivoting. Actually, it is possible to do better: This is the Cholesky factorization.

The following easy proposition shows that, in principle, A can be premultipled by some permutation matrix, P, so that PA can be converted to upper-triangular form without using any pivoting.

**Proposition 3.17** Let A be an invertible  $n \times n$ -matrix. Then, there is some permutation matrix, P, so that PA[1..k, 1..k] is invertible for k = 1, ..., n.

Proof. The case n=1 is trivial, and so is the case n=2 (we swap the rows if necessary). If  $n \geq 3$ , we proceed by induction. Since A is invertible, its columns are linearly independent; so, in particular, its first n-1 columns are also linearly independent. Delete the last column of A. Since the remaining n-1 columns are linearly independent, there are also n-1 linearly independent rows in the corresponding  $n \times (n-1)$  matrix. Thus, there is a permutation of these n rows so that the  $(n-1) \times (n-1)$  matrix consisting of the first n-1 rows is invertible. But, then, there is a corresponding permutation matrix,  $P_1$ , so that the first n-1 rows and columns of  $P_1A$  form an invertible matrix, A'. Applying the induction hypothesis to the  $(n-1) \times (n-1)$  matrix, A', we see that there some permutation matrix,  $P_2$ , so that  $P_2P_1A[1..k, 1..k]$  is invertible, for  $k=1,\ldots,n-1$ . Since A is invertible in the first place and  $P_1$  and  $P_2$  are invertible,  $P_1P_2A$  is also invertible, and we are done. □

**Remark:** One can also prove Proposition 3.17 using a clever reordering of the Gaussian elimination steps. Indeed, we know that if A is invertible, then there are permutation matrices,  $P_i$ , and products of elementary matrices,  $E_i$ , so that

$$A_n = E_{n-1}P_{n-1}\cdots E_2P_2E_1P_1A$$
,

where  $U = A_n$  is upper-triangular. For example, when n = 4, we have  $E_3P_3E_2P_2E_1P_1A = U$ . We can define new matrices  $E'_1, E'_2, E'_3$  which are still products of elementary matrices and we have

$$E_3' E_2' E_1' P_3 P_2 P_1 A = U.$$

Indeed, if we let  $E_3' = E_3$ ,  $E_2' = P_3 E_2 P_3^{-1}$ , and  $E_1' = P_3 P_2 E_1 P_2^{-1} P_3^{-1}$ , we easily verify that each  $E_k'$  is a product of elementary matrices and that

$$E_3'E_2'E_1'P_3P_2P_1 = E_3(P_3E_2P_3^{-1})(P_3P_2E_1P_2^{-1}P_3^{-1})P_3P_2P_1 = E_3P_3E_2P_2E_1P_1.$$

In general, we let

$$E'_{k} = P_{n-1} \cdots P_{k+1} E_{k} P_{k+1}^{-1} \cdots P_{n-1}^{-1},$$

and we have

$$E'_{n-1}\cdots E'_1P_{n-1}\cdots P_1A=U.$$

**Theorem 3.18** For every invertible  $n \times n$ -matrix, A, there is some permutation matrix, P, some upper-triangular matrix, U, and some unit lower-triangular matrix, L, so that PA = LU (recall,  $L_{ii} = 1$  for i = 1, ..., n). Furthermore, if P = I, then L and U are unique and they are produced as a result of Gaussian elimination without pivoting. Furthermore, if P = I, then L is simply obtained from the  $E_k^{-1}$ 's.

*Proof.* The only part that has not been proved is the uniqueness part and how L arises from the  $E_k^{-1}$ 's. Assume that A is invertible and that  $A = L_1U_1 = L_2U_2$ , with  $L_1, L_2$  unit lower-triangular and  $U_1, U_2$  upper-triangular. Then, we have

$$L_2^{-1}L_1 = U_2U_1^{-1}.$$

However, it is obvious that  $L_2^{-1}$  is lower-triangular and that  $U_1^{-1}$  is upper-triangular, and so,  $L_2^{-1}L_1$  is lower-triangular and  $U_2U_1^{-1}$  is upper-triangular. Since the diagonal entries of  $L_1$  and  $L_2$  are 1, the above equality is only possible if  $U_2U_1^{-1} = I$ , that is,  $U_1 = U_2$ , and so,  $L_1 = L_2$ . Finally, since  $L = E_1^{-1}E_2^{-1}\cdots E_{n-1}^{-1}$  and each  $E_k^{-1}$  is a product of elementary matrices of the form  $E_{i,k;-\beta_{i,k}}$  with  $k \leq i \leq n$ , we see that L is simply the lower-triangular matrix whose kth column is the kth column of  $E_k^{-1}$ , with  $1 \leq k \leq n-1$  (with  $L_{ii} = 1$  for  $i = 1, \ldots, n$ ).  $\square$ 

**Remark:** It can be shown that Gaussian elimination + back-substitution requires  $n^3/3 + O(n^2)$  additions,  $n^3/3 + O(n^2)$  multiplications and  $n^2/2 + O(n)$  divisions.

Let us now briefly comment on the choice of a pivot. Although theoretically, any pivot can be chosen, the possibility of roundoff errors implies that it is not a good idea to pick very small pivots. The following example illustrates this point. Consider the linear system

$$\begin{array}{rcl}
10^{-4}x & + & y & = & 1 \\
x & + & y & = & 2.
\end{array}$$

Since  $10^{-4}$  is nonzero, it can be taken as pivot, and we get

$$10^{-4}x + y = 1$$
$$(1 - 10^{4})y = 2 - 10^{4}.$$

Thus, the exact solution is

$$x = \frac{1}{1 - 10^{-4}}, \quad y = \frac{1 - 2 \times 10^{-4}}{1 - 10^{-4}}.$$

However, if roundoff takes place on the fourth digit, then  $1 - 10^4 = -9999$  and  $2 - 10^4 = -9998$  will be rounded off both to -9990, and then, the solution is x = 0 and y = 1, very far from the exact solution where  $x \approx 1$  and  $y \approx 1$ . The problem is that we picked a very small pivot. If instead we permute the equations, the pivot is 1, and after elimination, we get the system

$$x + y = 2$$
  
 $(1-10^{-4})y = 1-2 \times 10^{-4}.$ 

This time,  $1 - 10^{-4} = -0.9999$  and  $1 - 2 \times 10^{-4} = -0.9998$  are rounded off to 0.999 and the solution is x = 1, y = 1, much closer to the exact solution.

To remedy this problem, one may use the strategy of partial pivoting. This consists of choosing during step k  $(1 \le k \le n-1)$  one of the entries  $a_{ik}^k$  such that

$$|a_{i\,k}^k| = \max_{k \le p \le n} |a_{p\,k}^k|.$$

By maximizing the value of the pivot, we avoid dividing by undesirably small pivots.

**Remark:** A matrix, A, is called *strictly column diagonally dominant* iff

$$|a_{jj}| > \sum_{i=1, i \neq j}^{n} |a_{ij}|, \text{ for } j = 1, \dots, n$$

(resp. strictly row diagonally dominant iff

$$|a_{ii}| > \sum_{j=1, j \neq i}^{n} |a_{ij}|, \text{ for } i = 1, \dots, n.$$

It has been known for a long time (before 1900, say by Hadamard) that if a matrix, A, is strictly column diagonally dominant (resp. strictly row diagonally dominant), then it is invertible. (This is a good exercise, try it!) It can also be shown that if A is strictly column diagonally dominant, then Gaussian elimination with partial pivoting does not actually require pivoting.

Another strategy, called *complete pivoting*, consists in choosing some entry  $a_{ij}^k$ , where  $k \leq i, j \leq n$ , such that

$$|a_{ij}^k| = \max_{k \le p, q \le n} |a_{pq}^k|.$$

However, in this method, if the chosen pivot is not in column k, it is also necessary to permute columns. This is achieved by multiplying on the right by a permutation matrix.

However, complete pivoting tends to be too expansive in practice, and partial pivoting is the method of choice.

A special case where the LU-factorization is particularly efficient is the case of tridiagonal matrices, which we now consider.

Consider the tridiagonal matrix

$$A = \begin{pmatrix} b_1 & c_1 \\ a_2 & b_2 & c_2 \\ & a_3 & b_3 & c_3 \\ & \ddots & \ddots & \ddots \\ & & a_{n-2} & b_{n-2} & c_{n-2} \\ & & & a_{n-1} & b_{n-1} & c_{n-1} \\ & & & & a_n & b_n \end{pmatrix}.$$

Define the sequence

$$\delta_0 = 1$$
,  $\delta_1 = b_1$ ,  $\delta_k = b_k \delta_{k-1} - a_k c_{k-1} \delta_{k-2}$ ,  $2 \le k \le n$ .

**Proposition 3.19** If A is the tridiagonal matrix above, then  $\delta_k = \det(A[1..k, 1..k])$ , for  $k = 1, \ldots, n$ .

*Proof.* By expanding  $\det(A[1..k, 1..k])$  with respect to its last row, the proposition follows by induction on k.  $\square$ 

**Theorem 3.20** If A is the tridiagonal matrix above and  $\delta_k \neq 0$  for k = 1, ..., n, then A has the following LU-factorization:

$$A = \begin{pmatrix} 1 & & & & \\ a_2 \frac{\delta_0}{\delta_1} & 1 & & & \\ & a_3 \frac{\delta_1}{\delta_2} & 1 & & & \\ & & \ddots & \ddots & & \\ & & a_{n-1} \frac{\delta_{n-3}}{\delta_{n-2}} & 1 & & & \\ & & & a_n \frac{\delta_{n-2}}{\delta_{n-1}} & 1 \end{pmatrix} \begin{pmatrix} \frac{\delta_1}{\delta_0} & c_1 & & & \\ & \frac{\delta_2}{\delta_1} & c_2 & & & \\ & & \frac{\delta_3}{\delta_2} & c_3 & & \\ & & & \ddots & \ddots & \\ & & & \frac{\delta_{n-1}}{\delta_{n-2}} & c_{n-1} \\ & & & \frac{\delta_n}{\delta_{n-1}} \end{pmatrix}.$$

*Proof.* Since  $\delta_k = \det(A[1..k, 1..k]) \neq 0$  for k = 1, ..., n, by Theorem 3.18 (and Proposition 3.15), we know that A has a unique LU-factorization. Therefore, it suffices to check that

the proposed factorization works. We easily check that

$$(LU)_{k\,k+1} = c_k, \quad 1 \le k \le n-1$$

$$(LU)_{k\,k-1} = a_k, \quad 2 \le k \le n$$

$$(LU)_{k\,l} = 0, \quad |k-l| \ge 2$$

$$(LU)_{1\,1} = \frac{\delta_1}{\delta_0} = b_1$$

$$(LU)_{k\,k} = \frac{a_k c_{k-1} \delta_{k-2} + \delta_k}{\delta_{k-1}} = b_k, \quad 2 \le k \le n,$$

since  $\delta_k = b_k \delta_{k-1} - a_k c_{k-1} \delta_{k-2}$ .  $\square$ 

It follows that there is a simple method to solve a linear system, Ax = d, where A is tridiagonal (and  $\delta_k \neq 0$  for k = 1, ..., n). For this, it is convenient to "squeeze" the diagonal matrix,  $\Delta$ , defined such that  $\Delta_{kk} = \delta_k/\delta_{k-1}$ , into the factorization so that  $A = (L\Delta)(\Delta^{-1}U)$ , and if we let

$$z_1 = \frac{c_1}{b_1}, \quad z_k = c_k \frac{\delta_{k-1}}{\delta_k}, \quad 2 \le k \le n,$$

 $A = (L\Delta)(\Delta^{-1}U)$  is written as

$$A = \begin{pmatrix} \frac{c_1}{z_1} & & & & \\ a_2 & \frac{c_2}{z_2} & & & \\ & a_3 & \frac{c_3}{z_3} & & \\ & & \ddots & \ddots & \\ & & & a_{n-1} & \frac{c_{n-1}}{z_{n-1}} & & & \\ & & & & a_n & \frac{c_n}{z_n} \end{pmatrix} \begin{pmatrix} 1 & z_1 & & & \\ & 1 & z_2 & & & \\ & & 1 & z_3 & & \\ & & & \ddots & \ddots & \\ & & & 1 & z_{n-2} & \\ & & & & 1 & z_{n-1} & \\ & & & & & 1 & \end{pmatrix}$$

As a consequence, the system Ax = d can be solved by constructing three sequences: First, the sequence

$$z_1 = \frac{c_1}{b_1}, \quad z_k = \frac{c_k}{b_k - a_k z_{k-1}}, \quad k = 2, \dots, n,$$

corresponding to the recurrence  $\delta_k = b_k \delta_{k-1} - a_k c_{k-1} \delta_{k-2}$  and obtained by dividing both sides of this equation by  $\delta_{k-1}$ , next

$$w_1 = \frac{d_1}{b_1}$$
,  $w_k = \frac{d_k - a_k w_{k-1}}{b_k - a_k z_{k-1}}$ ,  $k = 2, \dots, n$ ,

corresponding to solving the system  $L\Delta w = d$ , and finally

$$x_n = w_n$$
,  $x_k = w_k - z_k x_{k+1}$ ,  $k = n - 1, n - 2, \dots, 1$ ,

corresponding to solving the system  $\Delta^{-1}Ux = w$ .

**Remark:** It can be verified that this requires 3(n-1) additions, 3(n-1) multiplications, and 2n divisions, a total of 8n-6 operations, which is much less that the  $O(2n^3/3)$  required by Gaussian elimination in general.

We now consider the special case of symmetric positive definite matrices. Recall that an  $n \times n$  symmetric matrix, A, is positive definite iff

$$x^{\top}Ax > 0$$
 for all  $x \in \mathbb{R}$  with  $x \neq 0$ .

Equivalently, A is symmetric positive definite iff all its eigenvalues are strictly positive. The following facts about a symmetric positive definite matrice, A, are easily established (some left as exercise):

- (1) The matrix A is invertible. (Indeed, if Ax = 0, then  $x^{T}Ax = 0$ , which implies x = 0.)
- (2) We have  $a_{ii} > 0$  for i = 1, ..., n. (Just observe that for  $x = e_i$ , the *i*th canonical basis vector of  $\mathbb{R}^n$ , we have  $e_i^{\top} A e_i = a_{ii} > 0$ .)
- (3) For every  $n \times n$  invertible matrix, Z, the matrix  $Z^{\top}AZ$  is symmetric positive definite iff A is symmetric positive definite.

Next, we prove that a symmetric positive definite matrix has a special LU-factorization of the form  $A = BB^{\top}$ , where B is a lower-triangular matrix whose diagonal elements are strictly positive. This is the *Cholesky factorization*.

First, we note that a symmetric positive definite matrix satisfies the condition of Proposition 3.15.

**Proposition 3.21** If A is a symmetric positive definite matrix, then A[1..k, 1..k] is invertible for k = 1, ..., n.

Proof. If  $w \in \mathbb{R}^k$ , with  $1 \le k \le n$ , we let  $x \in \mathbb{R}^n$  be the vector with  $x_i = w_i$  for  $i = 1, \ldots, k$  and  $x_i = 0$  for  $i = k + 1, \ldots, n$ . Now, since A is symmetric positive definite, we have  $x^{\top}Ax > 0$  for all  $x \in \mathbb{R}^n$  with  $x \ne 0$ . This holds in particular for all vectors x obtained from nonzero vectors  $x \in \mathbb{R}^n$  as defined earlier, which proves that each  $x \in \mathbb{R}^n$  is symmetric positive definite. Thus,  $x \in \mathbb{R}^n$  is also invertible.  $x \in \mathbb{R}^n$ 

Let A be a symmetric positive definite matrix and write

$$A = \begin{pmatrix} a_{11} & W^{\top} \\ W & B \end{pmatrix}.$$

Since A is symmetric positive definite,  $a_{11} > 0$ , and we can compute  $\alpha = \sqrt{a_{11}}$ . The trick is that we can factor A uniquely as

$$A = \begin{pmatrix} a_{11} & W^{\top} \\ W & B \end{pmatrix} = \begin{pmatrix} \alpha & 0 \\ W/\alpha & I \end{pmatrix} \begin{pmatrix} 1 & 0 \\ 0 & B - WW^{\top}/a_{11} \end{pmatrix} \begin{pmatrix} \alpha & W^{\top}/\alpha \\ 0 & I \end{pmatrix},$$

i.e., as  $A = B_1 A_1 B_1^{\top}$ , where  $B_1$  is lower-triangular with positive diagonal entries. Thus,  $B_1$  is invertible, and by fact (3) above,  $A_1$  is also symmetric positive definite.

**Theorem 3.22** (Cholesky Factorization) Let A be a symmetric positive definite matrix. Then, there is some lower-triangular matrix, B, so that  $A = BB^{\top}$ . Furthermore, B can be chosen so that its diagonal elements are strictly positive, in which case, B is unique.

*Proof.* We proceed by induction on k. For k=1, we must have  $a_{11}>0$ , and if we let  $\alpha=\sqrt{a_{11}}$  and  $B=(\alpha)$ , the theorem holds trivially. If  $k\geq 2$ , as we explained above, again we must have  $a_{11}>0$ , and we can write

$$A = \begin{pmatrix} a_{11} & W^{\top} \\ W & B \end{pmatrix} = \begin{pmatrix} \alpha & 0 \\ W/\alpha & I \end{pmatrix} \begin{pmatrix} 1 & 0 \\ 0 & B - WW^{\top}/a_{11} \end{pmatrix} \begin{pmatrix} \alpha & W^{\top}/\alpha \\ 0 & I \end{pmatrix} = B_1 A_1 B_1^{\top},$$

where  $\alpha = \sqrt{a_{11}}$ , the matrix  $B_1$  is invertible and

$$A_1 = \begin{pmatrix} 1 & 0 \\ 0 & B - WW^{\top}/a_{11} \end{pmatrix}$$

is symmetric positive definite. However, this implies that  $B - WW^{\top}/a_{11}$  is also symmetric positive definite (consider  $x^{\top}A_1x$  for every  $x \in \mathbb{R}^n$  with  $x \neq 0$  and  $x_1 = 0$ ). Thus, we can apply the induction hypothesis to  $B - WW^{\top}/a_{11}$ , and we find a unique lower-triangular matrix, L, with positive diagonal entries, so that

$$B - WW^{\top}/a_{11} = LL^{\top}.$$

But then, we get

$$A = \begin{pmatrix} \alpha & 0 \\ W/\alpha & I \end{pmatrix} \begin{pmatrix} 1 & 0 \\ 0 & B - WW^{\top}/a_{11} \end{pmatrix} \begin{pmatrix} \alpha & W^{\top}/\alpha \\ 0 & I \end{pmatrix}$$

$$= \begin{pmatrix} \alpha & 0 \\ W/\alpha & I \end{pmatrix} \begin{pmatrix} 1 & 0 \\ 0 & LL^{\top} \end{pmatrix} \begin{pmatrix} \alpha & W^{\top}/\alpha \\ 0 & I \end{pmatrix}$$

$$= \begin{pmatrix} \alpha & 0 \\ W/\alpha & I \end{pmatrix} \begin{pmatrix} 1 & 0 \\ 0 & L \end{pmatrix} \begin{pmatrix} 1 & 0 \\ 0 & L^{\top} \end{pmatrix} \begin{pmatrix} \alpha & W^{\top}/\alpha \\ 0 & I \end{pmatrix}$$

$$= \begin{pmatrix} \alpha & 0 \\ W/\alpha & L \end{pmatrix} \begin{pmatrix} \alpha & W^{\top}/\alpha \\ 0 & L^{\top} \end{pmatrix}.$$

Therefore, if we let

$$B = \begin{pmatrix} \alpha & 0 \\ W/\alpha & L \end{pmatrix},$$

we have a unique lower-triangular matrix with positive diagonal entries and  $A = BB^{\top}$ .  $\square$ 

**Remark:** If  $A = BB^{\top}$ , where B is any invertible matrix, then A is symmetric positive definite. Obviously,  $BB^{\top}$  is symmetric, and since B is invertible and

$$x^{\mathsf{T}}Ax = x^{\mathsf{T}}BB^{\mathsf{T}}x = (B^{\mathsf{T}}x)^{\mathsf{T}}B^{\mathsf{T}}x,$$

it is clear that  $x^{\top}Ax > 0$  if  $x \neq 0$ .

The proof of Theorem 3.22 immediately yields an algorithm to compute B from A. For  $j = 1, \ldots, n$ ,

$$b_{jj} = \left(a_{jj} - \sum_{k=1}^{j-1} b_{jk}^2\right)^{1/2},$$

and for  $i = j + 1, \ldots, n$ ,

$$b_{ij} = \left(a_{ij} - \sum_{k=1}^{j-1} b_{ik} b_{jk}\right) / b_{jj}.$$

The Cholseky factorization can be used to solve linear systems, Ax = b, where A is symmetric positive definite: Solve the two systems Bw = b and  $B^{\top}x = w$ .

**Remark:** It can be shown that this methods requires  $n^3/6 + O(n^2)$  additions,  $n^3/6 + O(n^2)$  multiplications,  $n^2/2 + O(n)$  divisions, and O(n) square root extractions. Thus, the Choleski method requires half of the number of operations required by Gaussian elimination (since Gaussian elimination requires  $n^3/3 + O(n^2)$  additions,  $n^3/3 + O(n^2)$  multiplications, and  $n^2/2 + O(n)$  divisions). It also requires half of the space (only B is needed, as opposed to both L and U). Furthermore, it can be shown that Cholesky's method is numerically stable.

For more on the stability analysis and efficient implementation methods of Gaussian elimination, LU-factoring and Cholesky factoring, see Demmel [14], Trefethen and Bau [52], Ciarlet [12], Golub and Van Loan [23], Strang [48, 49], and Kincaid and Cheney [28].

## 3.8 Futher Readings

Thorough expositions of the material covered in Chapter 2 and 3 can be found in Strang [49, 48], Lang [31], Artin [1], Mac Lane and Birkhoff [34], Bourbaki [7, 8], Van Der Waerden [53], Bertin [6], and Horn and Johnson [27]. These notions of linear algebra are nicely put to use in classical geometry, see Berger [3, 4], Tisseron [51] and Dieudonné [15].

Another rather complete reference is the text by Walter Noll [39]. But beware, the notation and terminology is a bit strange!

We now turn to polynomials.