Chapter 12 Solving Linear Equations

An Introduction to Optimization

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- Consider a system of linear equations Ax = b, where A ∈ R^{m×n} and b ∈ R^m, m ≥ n, and rank(A) = n. Note that the number of unknowns, n, is no larger than the number of equations, m.
- If b does not belong to the range of A, that is, b ∉ R(A), then this system of equations is said to be *inconsistent* or *overdetermined*.
- Our goal is to find the vector(s) x minimizing ||Ax b||². This problem is a special case of the nonlinear least-squares problem discussed in Section 9.4.

- Let x^* be a vector that minimizes $||Ax b||^2$; that is, for all $x \in \mathbb{R}^n$ $||Ax - b||^2 \ge ||Ax^* - b||^2$
- We refer to the vector x^* as a least-squares solution to Ax = b. In the case where Ax = b has a solution, then the solution is a least-squares solution. Otherwise, a least-squares solution minimizes the norm of the difference between the left- and right-hand sides of the equation Ax = b

- Lemma 12.1: Let $A \in \mathbb{R}^{m \times n}$, $m \ge n$. Then, rank(A) = n if and only if $rank(A^T A) = n$ (i.e., the square matrix $A^T A$ is nonsingular).
- Theorem 12.1: The unique vector x* that minimizes ||Ax b||² is given by the solution to the equation A^TAx = A^Tb; that is,

$$\boldsymbol{x}^* = (\boldsymbol{A}^T \boldsymbol{A})^{-1} \boldsymbol{A}^T \boldsymbol{b}$$

- The columns of A span the range R(A) of A, which is an n-dimensional subspace of R^m. The equation Ax = b has a solution if and only if b∈ R(A).
- If m = n, then $b \in \mathcal{R}(A)$ always, and the solution is $x^* = A^{-1}b$

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- Suppose now that m > n. We would expect the likelihood of b∈ R(A) to be small, because the subspace spanned by the columns of A is very "thin."
- Suppose that b ∉ R(A). We wish to find a point h ∈ R(A) that is "closest" to b. Geometrically, the point h should be such that the vector e = h − b is orthogonal to the subspace R(A)
- We call *h* the *orthogonal projection* of *b* onto the subspace $\mathcal{R}(A)$ It turns out that $h = Ax^* = A(A^TA)^{-1}A^Tb$.



- Write $A = [a_1, ..., a_n]$, where $a_1, ..., a_n$ are the columns of A
- ▶ The vector *e* is orthogonal to *R*(*A*) if and only if it is orthogonal to each of the columns *a*₁,..., *a*_n of *A*.
- Note that ⟨e, a_i⟩ = 0, i = 1, ..., n if and only if for any set of scalars {x₁, x₂, ..., x_n}, we also have ⟨e, x₁a₁ + ··· + x_na_n⟩ = 0

Any vector in $\mathcal{R}(\mathbf{A})$ has the form $x_1\mathbf{a}_1 + \cdots + x_n\mathbf{a}_n$

- Proposition 12.1: Let $h \in \mathcal{R}(A)$ be such that h b is orthogonal to $\mathcal{R}(A)$. Then, $h = Ax^* = A(A^TA)^{-1}A^Tb$
- Note that the matrix

$$oldsymbol{A}^Toldsymbol{A} = egin{bmatrix} \langleoldsymbol{a}_1,oldsymbol{a}_1
angle & \cdots & \langleoldsymbol{a}_n,oldsymbol{a}_1
angle \ dots\ egin{matrix} oldsymbol{a}_n,oldsymbol{a}_n
angle \ dots\ eta_n,oldsymbol{a}_n
angle \ eta_n,oldsymbol{a}_n
angle \ \cdots & \langleoldsymbol{a}_n,oldsymbol{a}_n
angle \end{bmatrix}$$

plays an important role in the least-squares solution. This matrix is often called the *Gram matrix* (or *Grammian*).

- Suppose that you are given two different types of concrete. The first type contains 30% cement, 40% gravel, and 30% sand (all percentages of weight). The second type contains 10% cement, 20% gravel, and 70% sand. How many pounds of each type of concrete should you mix together so that you get a concrete mixture that has as close as possible to a total of 5 pounds of cement, 3 pounds of gravel, and 4 pounds of sand?
- The problem can be formulated as a least-squares problem with

$$\boldsymbol{A} = \begin{bmatrix} 0.3 & 0.1 \\ 0.4 & 0.2 \\ 0.3 & 0.7 \end{bmatrix} \qquad \boldsymbol{b} = \begin{bmatrix} 5 \\ 3 \\ 4 \end{bmatrix}$$

where the decision variable is $x = [x_1, x_2]^T$ and x_1 and x_2 are the amounts of concrete of the first and second types, respectively.

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• After some algebra, we obtain the solution:

$$\boldsymbol{x}^* = (\boldsymbol{A}^T \boldsymbol{A})^{-1} \boldsymbol{A}^T \boldsymbol{b}$$

= $\frac{1}{(0.34)(0.54) - (0.32)^2} \begin{bmatrix} 0.54 & -0.32 \\ -0.32 & 0.34 \end{bmatrix} \begin{bmatrix} 3.9 \\ 3.9 \end{bmatrix} = \begin{bmatrix} 10.6 \\ 0.961 \end{bmatrix}$

 Line Fitting. Suppose that a process has a single input t∈ R and a single output y∈ R. Suppose that we perform an experiment on the process, resulting in a number of measurements.

i	0	1	2
t_i	2	3	4
y_i	3	4	15

The *i*th measurement results in the input labeled t_i and the output labeled y_i . We would like to find a straight line given by y = mt + c that fits the experimental data. In other words, we wish to find two numbers, m and c, such that $y_i = mt_i + c$, i = 0, 1, 2.

• However, it is apparent that there is no choice of *m* and *c* that results in the requirement above. Therefore, we would like to find the values of *m* and *c* that best fit the data.



• We can represent our problem as a system of linear equations of the form 2m + c = 3

$$3m + c = 4$$

$$4m + c = 15$$

$$\boldsymbol{A} = \begin{bmatrix} 2 & 1 \\ 3 & 1 \\ 4 & 1 \end{bmatrix} \quad \boldsymbol{b} = \begin{bmatrix} 3 \\ 4 \\ 15 \end{bmatrix} \quad \boldsymbol{x} = \begin{bmatrix} m \\ c \end{bmatrix}$$

Notice that since rank(A) < rank([A, b]), the vector **b** does not belong to the range of **A**. Thus, this system is inconsistent.

> The solution to this least-squares problem is

$$\boldsymbol{x}^* = \begin{bmatrix} m^* \\ c^* \end{bmatrix} = (\boldsymbol{A}^T \boldsymbol{A})^{-1} \boldsymbol{A}^T \boldsymbol{b} = \begin{bmatrix} 6 \\ -32/3 \end{bmatrix}$$

- Attenuation Estimation. A wireless transmitter sends a discrete-time signal {s₀, s₁, s₂} (of duration 3) to a receiver. The real number s_i is the value of the signal at time i.
- The transmitted signal takes two paths to the receiver: a direct path, with delay 10 and attenuation factor a₁, and an indirect (reflected) path, with delay 12 and attenuation factor a₂. The received signal is the sum of the signals from these two paths, with their respective delays and attenuation factors.



Suppose that the received signal is measured from times 10 through 14 as r₁₀, r₁₁, ..., r₁₄. We wish to compute the least-squares estimates of a₁ and a₂, based on the following values

s_0	s_1	s_2	r_{10}	r_{11}	r_{12}	r_{13}	r_{14}
1	2	1	4	7	8	6	3

> The problem can be posed as a least-squares problem with

$$a_{1}s + a_{2}s' = r$$

$$a_{1}s_{0} + 0 = r_{10}$$

$$a_{1}s_{1} + 0 = r_{11}$$

$$a_{1}s_{2} + a_{2}s_{0} = r_{12}$$

$$0 + a_{2}s_{1} = r_{13}$$

$$0 + a_{2}s_{2} = r_{14}$$

$$A = \begin{bmatrix} s_{0} & 0 \\ s_{1} & 0 \\ s_{2} & s_{0} \\ 0 & s_{1} \\ 0 & s_{2} \end{bmatrix}$$

$$x = \begin{bmatrix} a_{1} \\ a_{2} \end{bmatrix}$$

$$b = \begin{bmatrix} r_{10} \\ r_{11} \\ r_{12} \\ r_{13} \\ r_{14} \end{bmatrix}$$

The least-squares estimate is given

$$\boldsymbol{x}^{*} = \begin{bmatrix} a_{1}^{*} \\ a_{2}^{*} \end{bmatrix} = (\boldsymbol{A}^{T}\boldsymbol{A})^{-1}\boldsymbol{A}^{T}\boldsymbol{b}$$
$$= \begin{bmatrix} \|\boldsymbol{s}\|^{2} & s_{0}s_{2} \\ s_{0}s_{2} & \|\boldsymbol{s}\|^{2} \end{bmatrix}^{-1} \begin{bmatrix} s_{0}r_{10} + s_{1}r_{11} + s_{2}r_{12} \\ s_{0}r_{12} + s_{1}r_{13} + s_{2}r_{14} \end{bmatrix}$$
$$= \begin{bmatrix} 6 & 1 \\ 1 & 6 \end{bmatrix}^{-1} \begin{bmatrix} 4 + 14 + 8 \\ 8 + 12 + 3 \end{bmatrix}$$
$$= \frac{1}{35} \begin{bmatrix} 133 \\ 112 \end{bmatrix}$$

• *Discrete Fourier Series*. Suppose that we are given a discrete time signal, represented by the vector

$$\boldsymbol{b} = [b_1, b_2, ..., b_m]^T$$

We wish to approximate this signal by a sum of sinusoids. Specifically, we approximate *b* by the vector

$$y_0 \boldsymbol{c}^{(0)} + \sum_{k=1}^n \left(y_k \boldsymbol{c}^{(k)} + z_k \boldsymbol{s}^{(k)} \right)$$

where $y_0, y_1, ..., y_n, z_1, ..., z_n \in R$ and the vectors $c^{(k)}$ and $s^{(k)}$ are given by

$$\boldsymbol{c}^{(0)} = \begin{bmatrix} \frac{1}{\sqrt{2}}, \frac{1}{\sqrt{2}}, \dots, \frac{1}{\sqrt{2}} \end{bmatrix}^{T}$$
$$\boldsymbol{c}^{(k)} = \begin{bmatrix} \cos\left(1\frac{2k\pi}{m}\right), \cos\left(2\frac{2k\pi}{m}\right), \dots, \cos\left(m\frac{2k\pi}{m}\right) \end{bmatrix}^{T}, k = 1, \dots, n$$
$$\boldsymbol{s}^{(k)} = \begin{bmatrix} \sin\left(1\frac{2k\pi}{m}\right), \sin\left(2\frac{2k\pi}{m}\right), \dots, \sin\left(m\frac{2k\pi}{m}\right) \end{bmatrix}^{T}, k = 1, \dots, n$$

• We call the sum of sinusoids above a *discrete Fourier series*. We wish to find $y_0, y_1, ..., y_n, z_1, ..., z_n$ such that $\left\| \left(y_0 \boldsymbol{c}^{(0)} + \sum_{k=1}^n \left(y_k \boldsymbol{c}^{(k)} + z_k \boldsymbol{s}^{(k)} \right) \right) - \boldsymbol{b} \right\|^2$

is minimized.

• To proceed, we define $A = [c^{(0)}, c^{(1)}, ..., c^{(n)}, s^{(1)}, ..., s^{(n)}]$ $x = [y_0, y_1, ..., y_n, z_1, ..., z_n]^T$ Our problem can be reformulated as minimizing $\|Ax - b\|^2$

We assume that m ≥ 2n + 1. To find the solution, we first compute A^TA. We make use of the following trigonometric identities: for any nonzero integer k that is not an integral multiple of m, we have

$$\sum_{i=1}^{m} \cos\left(i\frac{2k\pi}{m}\right) = 0$$
$$\sum_{i=1}^{m} \sin\left(i\frac{2k\pi}{m}\right) = 0$$

with the aid of these identities, we can verify that

$$\boldsymbol{c}^{(k)T}\boldsymbol{c}^{(j)} = \begin{cases} m/2 & \text{if } k = j \\ 0 & \text{otherwise} \end{cases}$$
$$\boldsymbol{s}^{(k)T}\boldsymbol{s}^{(j)} = \begin{cases} m/2 & \text{if } k = j \\ 0 & \text{otherwise} \end{cases}$$

 $\boldsymbol{c}^{(k)T}\boldsymbol{s}^{(j)} = 0$ for any k, j

• Hence, $\mathbf{A}^T \mathbf{A} = \frac{m}{2} \mathbf{I}_{2n+1}$, which is clearly nonsingular, with inverse $(\mathbf{A}^T \mathbf{A})^{-1} = \frac{2}{m} \mathbf{I}_{2n+1}$ Therefore, the solution to our problem is $\mathbf{x}^* = [y_0^*, y_1^*, ..., y_n^*, z_1^*, ..., z_n^*]^T$ $= (\mathbf{A}^T \mathbf{A})^{-1} \mathbf{A}^T \mathbf{b}$ $= \frac{2}{m} \mathbf{A}^T \mathbf{b}$

We represent the solution as

$$y_0^* = \frac{\sqrt{2}}{m} \sum_{i=1}^m b_i \qquad y_k^* = \frac{2}{m} \sum_{i=1}^m b_i \cos\left(i\frac{2k\pi}{m}\right) \qquad k = 1, ..., n$$
$$z_k^* = \frac{2}{m} \sum_{i=1}^m b_i \sin\left(i\frac{2k\pi}{m}\right) \qquad k = 1, ..., n$$

We call these *discrete Fourier coefficients*.

Orthogonal Projections. Let V ⊂ Rⁿ be a subspace. Given a vector x ∈ Rⁿ, we write the orthogonal decomposition of x as x = x_V + x_{V[⊥]}

where $x_{\mathcal{V}} \in \mathcal{V}$ is the orthogonal projection of x onto \mathcal{V} and $x_{\mathcal{V}^{\perp}} \in \mathcal{V}^{\perp}$ is the orthogonal projection of x onto \mathcal{V}^{\perp} . We can write $x_{\mathcal{V}} = Px$ for some matrix P called *orthogonal projector*. In the following, we derive expressions for P for the case where $\mathcal{V} = \mathcal{R}(A)$ and the case where $\mathcal{V} = \mathcal{N}(A)$

Consider a matrix A ∈ R^{m×n}, m ≥ n, and rank(A) = n. Let V = R(A) be the range of A. In this case we can write an expression for P in terms of A

- By Proposition 12.1 we have $\boldsymbol{x}_{\mathcal{V}} = \boldsymbol{A}(\boldsymbol{A}^T\boldsymbol{A})^{-1}\boldsymbol{A}^T\boldsymbol{x}$, whence $\boldsymbol{P} = \boldsymbol{A}(\boldsymbol{A}^T\boldsymbol{A})^{-1}\boldsymbol{A}^T$. Note that by Proposition 12.1, we may also write $\boldsymbol{x}_{\mathcal{V}} = \arg\min_{\boldsymbol{y}\in\mathcal{V}} ||\boldsymbol{y} - \boldsymbol{x}||$
- Next, consider a matrix A ∈ R^{m×n}, m ≤ n, and rank(A) = m. Let V = N(A) be the nullspace of A. To derive an expression for the orthogonal projector P in terms of A for this case, we use the formula derived above and the identity N(A)[⊥] = R(A^T) (see Theorem 3.4).

Indeed, if $\mathcal{U} = \mathcal{R}(\mathbf{A}^T)$, then the orthogonal decomposition with respect to \mathcal{U} is $\mathbf{x} = \mathbf{x}_{\mathcal{U}} + \mathbf{x}_{\mathcal{U}^{\perp}}$, where $\mathbf{x}_{\mathcal{U}} = \mathbf{A}^T (\mathbf{A}\mathbf{A}^T)^{-1} \mathbf{A}\mathbf{x}$

• Because $\mathcal{N}(\mathbf{A})^{\perp} = \mathcal{R}(\mathbf{A}^T)$, we deduce that $\mathbf{x}_{\mathcal{V}^{\perp}} = \mathbf{x}_{\mathcal{U}} = \mathbf{A}^T (\mathbf{A}\mathbf{A}^T)^{-1} \mathbf{A}\mathbf{x}$ Hence,

$$\boldsymbol{x}_{\mathcal{V}} = \boldsymbol{x} - \boldsymbol{x}_{\mathcal{V}^{\perp}} = \boldsymbol{x} - \boldsymbol{A}^{T}(\boldsymbol{A}\boldsymbol{A}^{T})^{-1}\boldsymbol{A}\boldsymbol{x} = (\boldsymbol{I} - \boldsymbol{A}^{T}(\boldsymbol{A}\boldsymbol{A}^{T})^{-1}\boldsymbol{A})\boldsymbol{x}$$

Thus, the orthogonal projector in this case is
 $\boldsymbol{P} = \boldsymbol{I} - \boldsymbol{A}^{T}(\boldsymbol{A}\boldsymbol{A}^{T})^{-1}\boldsymbol{A}$

Assume that we are originally given three experimental results (t₀, y₀), (t₁, y₁), (t₂, y₂), and we find the parameters m* and c* of the straight line that best bits these data. Suppose that we are now given an extra measurement point (t₃, y₃). We can use previous calculations of m* and c* for the three data points to calculate the parameters for the four data points. This procedure is called the *recursive least-squares* (RLS) algorithm.

• Consider the problem of minimizing $\|\boldsymbol{A}_0\boldsymbol{x} - \boldsymbol{b}^{(0)}\|^2$. The solution is given by $\boldsymbol{x}^{(0)} = \boldsymbol{G}_0^{-1}\boldsymbol{A}_0^T\boldsymbol{b}^{(0)}$, where $\boldsymbol{G}_0 = \boldsymbol{A}_0^T\boldsymbol{A}_0$ Consider now the problem of minimizing $\|\boldsymbol{\Gamma}_{\boldsymbol{A}}\|^2 = \|\boldsymbol{\Gamma}_{\boldsymbol{b}}\|^2$

$$\left\| \begin{bmatrix} oldsymbol{A}_0 \ oldsymbol{A}_1 \end{bmatrix} oldsymbol{x} - \begin{bmatrix} oldsymbol{b}^{(0)} \ oldsymbol{b}^{(1)} \end{bmatrix}
ight\|$$
iven by

The solution is given by

$$oldsymbol{x}^{(1)} = oldsymbol{G}_1^{-1} egin{bmatrix} oldsymbol{A}_0 \ oldsymbol{A}_1 \end{bmatrix}^T egin{bmatrix} oldsymbol{b}^{(0)} \ oldsymbol{b}^{(1)} \end{bmatrix} \qquad oldsymbol{G}_1 = egin{bmatrix} oldsymbol{A}_0 \ oldsymbol{A}_1 \end{bmatrix}^T egin{bmatrix} oldsymbol{A}_0 \ oldsymbol{A}_1 \end{bmatrix}^T egin{bmatrix} oldsymbol{A}_0 \ oldsymbol{A}_1 \end{bmatrix}^T egin{bmatrix} oldsymbol{A}_0 \ oldsymbol{A}_1 \end{bmatrix}$$

Our goal is to write $x^{(1)}$ as a function of $x^{(0)}$, G_0 , and the new data A_1 and $b^{(1)}$

• To this end, we first write G_1 as

$$oldsymbol{G}_1 = [oldsymbol{A}_0^T oldsymbol{A}_1^T] egin{bmatrix} oldsymbol{A}_0 \ oldsymbol{A}_1 \end{bmatrix} = oldsymbol{A}_0^T oldsymbol{A}_0 + oldsymbol{A}_1^T oldsymbol{A}_1 = oldsymbol{G}_0 + oldsymbol{A}_1^T oldsymbol{A}_1$$

Next, we write

$$\begin{bmatrix} \boldsymbol{A}_0 \\ \boldsymbol{A}_1 \end{bmatrix}^T \begin{bmatrix} \boldsymbol{b}^{(0)} \\ \boldsymbol{b}^{(1)} \end{bmatrix} = \begin{bmatrix} \boldsymbol{A}_0^T & \boldsymbol{A}_1^T \end{bmatrix} \begin{bmatrix} \boldsymbol{b}^{(0)} \\ \boldsymbol{b}^{(1)} \end{bmatrix} = \boldsymbol{A}_0^T \boldsymbol{b}^{(0)} + \boldsymbol{A}_1^T \boldsymbol{b}^{(1)}$$

To proceed further, we write $A_0^T b^{(0)}$ as

$$A_0^T b^{(0)} = G_0 G_0^{-1} A_0^T b^{(0)} = G_0 x^{(0)}$$

= $(G_1 - A_1^T A_1) x^{(0)} = G_1 x^{(0)} - A_1^T A_1 x^{(0)}$

• Combining these formulas, we see that we can write $x^{(1)}$ as

$$\begin{aligned} \boldsymbol{x}^{(1)} &= \boldsymbol{G}_{1}^{-1} \begin{bmatrix} \boldsymbol{A}_{0} \\ \boldsymbol{A}_{1} \end{bmatrix}^{T} \begin{bmatrix} \boldsymbol{b}^{(0)} \\ \boldsymbol{b}^{(1)} \end{bmatrix} = \boldsymbol{G}_{1}^{-1} \Big(\boldsymbol{G}_{1} \boldsymbol{x}^{(0)} - \boldsymbol{A}_{1}^{T} \boldsymbol{A}_{1} \boldsymbol{x}^{(0)} + \boldsymbol{A}_{1}^{T} \boldsymbol{b}^{(1)} \Big) \\ &= \boldsymbol{x}^{(0)} + \boldsymbol{G}_{1}^{-1} \boldsymbol{A}_{1}^{T} \Big(\boldsymbol{b}^{(1)} - \boldsymbol{A}_{1} \boldsymbol{x}^{(0)} \Big) \end{aligned}$$

where G_1 can be calculated using $G_1 = G_0 + A_1^T A_1$

With this formula, x⁽¹⁾ can be computed using only x⁽⁰⁾, A₁, b⁽¹⁾ and G₀. Hence, we have a way of using our previous efforts in calculating x⁽⁰⁾ to compute x⁽¹⁾. Observe that if the new data are consistent with the old data, that is, A₁x⁽⁰⁾ = b⁽¹⁾, then the correction term is 0 and the updated solution x⁽¹⁾ is equal to the previous solution x⁽⁰⁾.

• At the (k + 1) th iteration, we have

$$m{G}_{k+1} = m{G}_k + m{A}_{k+1}^T m{A}_{k+1}$$

 $m{x}^{(k+1)} = m{x}^{(k)} + m{G}_{k+1}^{-1} m{A}_{k+1}^T \Big(m{b}^{(k+1)} - m{A}_{k+1} m{x}^{(k)} \Big)$

The vector $\boldsymbol{b}^{(k+1)} - \boldsymbol{A}_{k+1}\boldsymbol{x}^{(k)}$ is often called the *innovation*. As before, observe that if the innovation is zero, then the updated solution $\boldsymbol{x}^{(k+1)}$ is equal to the previous solution $\boldsymbol{x}^{(k)}$

We can see that to compute x^(k+1) we need G⁻¹_{k+1} rather than
 G_{k+1}. It turns out that we can derive an update formula for G⁻¹_{k+1} itself.

• Lemma 12.2: Let A be a nonsingular matrix. Let U and V be matrices such that $I + VA^{-1}U$ is nonsingular. Then, A + UV is nonsingular, and

$$(A + UV)^{-1} = A^{-1} - (A^{-1}U)(I + VA^{-1}U)^{-1}(VA^{-1})$$

• By Lemma 12.2, we get

$$G_{k+1}^{-1} = (G_k + A_{k+1}^T A_{k+1})^{-1}$$

 $= G_k^{-1} - G_k^{-1} A_{k+1}^T (I + A_{k+1} G_k^{-1} A_{k+1}^T)^{-1} A_{k+1} G_k^{-1}$

For simplicity of notation, we rewrite G_k^{-1} as P_k . We summarize by writing the RLS algorithm using P_k

$$\boldsymbol{P}_{k+1} = \boldsymbol{P}_k - \boldsymbol{P}_k \boldsymbol{A}_{k+1}^T (\boldsymbol{I} + \boldsymbol{A}_{k+1} \boldsymbol{P}_k \boldsymbol{A}_{k+1}^T)^{-1} \boldsymbol{A}_{k+1} \boldsymbol{P}_k$$
$$\boldsymbol{x}^{(k+1)} = \boldsymbol{x}^{(k)} + \boldsymbol{P}_{k+1} \boldsymbol{A}_{k+1}^T \left(\boldsymbol{b}^{(k+1)} - \boldsymbol{A}_{k+1} \boldsymbol{x}^{(k)} \right)$$

 In the special case where the new data at each step are such that *A*_{k+1} is a matrix consisting of a single row, *A*_{k+1} = *a*^T_{k+1}, and *b*^(k+1) is a scalar, *b*^(k+1) = *b*_{k+1}, we get

$$oldsymbol{P}_{k+1} = oldsymbol{P}_k - rac{oldsymbol{P}_k oldsymbol{a}_{k+1} oldsymbol{a}_{k+1}^T oldsymbol{P}_k}{1 + oldsymbol{a}_{k+1}^T oldsymbol{P}_k oldsymbol{a}_{k+1}}
onumber \ oldsymbol{x}^{(k+1)} = oldsymbol{x}^{(k)} + oldsymbol{P}_{k+1} oldsymbol{a}_{k+1} ildsymbol{b}_{k+1} - oldsymbol{a}_{k+1}^T oldsymbol{x}^{(k)} ildsymbol{)}$$

Example
• Let
$$A_{0} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \\ 1 & 1 \end{bmatrix} \quad b^{(0)} = \begin{bmatrix} 1 \\ 1 \\ 1 \end{bmatrix}$$

$$A_{1} = a_{1}^{T} = \begin{bmatrix} 2 & 1 \end{bmatrix} \quad b^{(1)} = b_{1} = \begin{bmatrix} 3 \end{bmatrix}$$

$$A_{2} = a_{2}^{T} = \begin{bmatrix} 3 & 1 \end{bmatrix} \quad b^{(2)} = b_{2} = \begin{bmatrix} 4 \end{bmatrix}$$

First compute the vector $\mathbf{x}^{(0)}$ minimizing $\|\mathbf{A}_0\mathbf{x} - \mathbf{b}^{(0)}\|^2$. Then, use the RLS algorithm to find $\mathbf{x}^{(2)}$ minimizing

$$\begin{aligned} \left\| \begin{bmatrix} \boldsymbol{A}_0 \\ \boldsymbol{A}_1 \\ \boldsymbol{A}_2 \end{bmatrix} \boldsymbol{x} - \begin{bmatrix} \boldsymbol{b}^{(0)} \\ \boldsymbol{b}^{(1)} \\ \boldsymbol{b}^{(2)} \end{bmatrix} \right\|^2 \end{aligned}$$

We have
$$\boldsymbol{P}_0 = (\boldsymbol{A}_0^T \boldsymbol{A}_0)^{-1} = \begin{bmatrix} 2/3 & -1/3 \\ -1/3 & 2/3 \end{bmatrix}$$
$$\boldsymbol{x}^{(0)} = \boldsymbol{P}_0 \boldsymbol{A}_0^T \boldsymbol{b}^{(0)} = \begin{bmatrix} 2/3 \\ 2/3 \end{bmatrix}$$

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• Applying the RLS algorithm twice, we get

$$P_{1} = P_{0} - \frac{P_{0}a_{1}a_{1}^{T}P_{0}}{1 + a_{1}^{T}P_{0}a_{1}} = \begin{bmatrix} 1/3 & -1/3 \\ -1/3 & 2/3 \end{bmatrix}$$
$$x^{(1)} = x^{(0)} + P_{1}a_{1}\left(b_{1} - a_{1}^{T}x^{(0)}\right) = \begin{bmatrix} 1 \\ 2/3 \end{bmatrix}$$
$$P_{2} = P_{1} - \frac{P_{1}a_{2}a_{2}^{T}P_{1}}{1 + a_{2}^{T}P_{1}a_{2}} = \begin{bmatrix} 1/6 & -1/4 \\ -1/4 & 5/8 \end{bmatrix}$$
$$x^{(2)} = x^{(1)} + P_{2}a_{2}\left(b_{2} - a_{2}^{T}x^{(1)}\right) = \begin{bmatrix} 13/12 \\ 5/8 \end{bmatrix}$$

Solution to A Linear Equation with Minimum Norm

- Consider now a system of linear equations Ax = b, where A ∈ R^{m×n}, b ∈ R^m, m ≤ n, and rank(A) = m. Note that the number of equations is no longer than the number of unknowns. There may exist an infinite number of solutions to this system of equations.
- However, as we shall see, there is only one solution that is closest to the origin: the solution to Ax = b whose norm ||x| is minimal.
- Let x^* be this solution that is $Ax^* = b$ and $||x^*|| \le ||x|$ for any x such that Ax = b. In other words, x^* is the solution to the problem

minimize $\|\boldsymbol{x}\|$

subject to Ax = b

Solution to A Linear Equation with Minimum Norm

- Theorem 12.2: The unique solution x* to Ax = b that minimizes the norm ||x| is given by
 x* = A^T(AA^T)⁻¹b
- Example: Find the point closest to the origin of R³ on the line of intersection of the two planes defined by the following two equations: x₁ + 2x₂ x₃ = 1

$$4x_1 + x_2 + 3x_3 = 0$$

Note that this problem is equivalent to the problem

minimize
$$\|\boldsymbol{x}\|$$

subject to $\boldsymbol{A}\boldsymbol{x} = \boldsymbol{b}$ $\boldsymbol{A} = \begin{bmatrix} 1 & 2 & -1 \\ 4 & 1 & 3 \end{bmatrix}$ $\boldsymbol{b} = \begin{bmatrix} 1 \\ 0 \end{bmatrix}$
 $\boldsymbol{x}^* = \boldsymbol{A}^T (\boldsymbol{A}\boldsymbol{A}^T)^{-1} \boldsymbol{b} = \begin{bmatrix} 0.0952 \\ 0.3333 \\ -0.2381 \end{bmatrix}$

- Kaczmarz's algorithm converges to the vector x* = A^T(AA^T)⁻¹b without explicitly having to invert the matrix AA^T. This is important from a practical point of view, especially when A has many rows.
- Let a_j^T denote the *j*th row of A, and b_j the *j*th component of band μ a positive scalar, $0 < \mu < 2$. Kaczmarz's algorithm is:

▶ 1. Set
$$i := 0$$
, initial condition $\boldsymbol{x}^{(0)}$

> 2. For
$$j = 1, ..., m$$
, set
 $\boldsymbol{x}^{(im+j)} = \boldsymbol{x}^{(im+j-1)} + \mu(b_j - \boldsymbol{a}_j^T \boldsymbol{x}^{(im+j-1)}) \frac{\boldsymbol{a}_j}{\boldsymbol{a}_j^T \boldsymbol{a}_j}$

• 3. Set i := i + 1; go to step 2.

• For the first m iterations, we have

$$\boldsymbol{x}^{(k+1)} = \boldsymbol{x}^{(k)} + \mu (b_{k+1} - \boldsymbol{a}_{k+1}^T \boldsymbol{x}^{(k)}) \frac{\boldsymbol{a}_{k+1}}{\boldsymbol{a}_{k+1}^T \boldsymbol{a}_{k+1}}$$

where, in each iteration, we use rows of A and corresponding components of b successively. For the (m + 1) th iteration, we revert back to the first row of A and the first component of b; that is, $x^{(m+1)} = x^{(m)} + \mu(b_1 - a_1^T x^{(m)}) \frac{a_1}{a_1^T a_1}$

We continue with the (m+2) th iteration using the second row of A and the second component of b, and so on, repeating the cycle every m iterations. The reason for $0 < \mu < 2$ will become apparent from the convergence anlaysis.

• Theorem 12.3: In Kaczmarz's algorithm, if $\boldsymbol{x}^{(0)} = \boldsymbol{0}$, then $\boldsymbol{x}^{(k)} \rightarrow \boldsymbol{x}^* = \boldsymbol{A}^T (\boldsymbol{A} \boldsymbol{A}^T)^{-1} \boldsymbol{b}$ as $k \rightarrow \infty$.

• Example: Let

$$\boldsymbol{A} = \begin{bmatrix} 1 & -1 \\ 0 & 1 \end{bmatrix} \qquad \boldsymbol{b} = \begin{bmatrix} 2 \\ 3 \end{bmatrix}$$

In this case $\mathbf{x}^* = [5,3]^T$. This figure shows a few iterations of Kaczmarz's algorithm with $\mu = 1$ and $\mathbf{x}^{(0)} = \mathbf{0}$. We have $\mathbf{a}_1^T = [1,-1]$ $\mathbf{a}_2^T = [0,1]$, $b_1 = 2$, $b_2 = 3$. The diagonal line passing through the point $[2,0]^T$ corresponds to the set $\{\mathbf{x} : \mathbf{a}_1^T \mathbf{x} = b_1\}$, and the horizontal line passing through the point $[0,3]^T$ corresponds to the set $\{\mathbf{x} : \mathbf{a}_2^T \mathbf{x} = b_2\}$.

• We perform three iterations:

$$\boldsymbol{x}^{(1)} = \begin{bmatrix} 0\\0 \end{bmatrix} + (2-0)^{\frac{1}{2}} \begin{bmatrix} 1\\-1 \end{bmatrix} = \begin{bmatrix} 1\\-1 \end{bmatrix}$$
$$\boldsymbol{x}^{(2)} = \begin{bmatrix} 1\\-1 \end{bmatrix} + (3-(-1)) \begin{bmatrix} 0\\1 \end{bmatrix} = \begin{bmatrix} 1\\3 \end{bmatrix}$$
$$\boldsymbol{x}^{(3)} = \begin{bmatrix} 1\\3 \end{bmatrix} + (2-(-2))^{\frac{1}{2}} \begin{bmatrix} 1\\-1 \end{bmatrix} = \begin{bmatrix} 3\\1 \end{bmatrix}$$



- Consider a system of linear equations Ax = b, where A ∈ R^{m×n} and rank(A) = r. Note that we always have r ≤ min{m, n}. In the case A ∈ R^{n×n} and rank(A) = n, the unique solution to the equation above has the form x* = A⁻¹b. Thus, to solve the problem in this case it is enough to know the inverse A⁻¹.
- A general approach to solving Ax = b. The approach involves defining a *pseudoinverse* or *generalized inverse* of a given matrix A ∈ R^{m×n}, which plays the role of A⁻¹ when A does not have an inverse. In particular, we discuss the *Moore-Penrose inverse* of a given matrix A, denoted A[†].

- Lemma 12.3 *Full-Rank Factorization*: Let $A \in \mathbb{R}^{m \times n}$, $rank(A) = r \le \min\{m, n\}$. Then, there exist matrices $B \in \mathbb{R}^{m \times r}$ and $C \in \mathbb{R}^{r \times n}$ such that A = BC, where rank(A) = rank(B) = rank(C) = r
- Proof: Because rank(A) = r, we can find r linearly independent columns of A. Without loss of generality, let $a_1, a_2, ..., a_r$ be such columns, where a_i is the i th column of A. The remaining columns of A can be expressed as linear combination of $a_1, a_2, ..., a_r$. Thus, a possible choice for B and C are $B = [a_1, a_2, ..., a_r] \in \mathbb{R}^{m \times r}$ $C = \begin{bmatrix} 1 \cdots 0 & c_{1,r+1} \cdots & c_{1,n} \\ \vdots & \ddots & \vdots & \vdots & \ddots & \vdots \\ 0 \cdots & 1 & c_{r,r+1} \cdots & c_{r,n} \end{bmatrix} \in \mathbb{R}^{r \times n}$

where the entries $c_{i,j}$ are such that for each j = r + 1, ..., n, we have $a_j = c_{1,j}a_1 + \cdots + c_{r,j}a_r$. Thus, A = BC

- Note that if m < n and $rank(\mathbf{A}) = m$, then we take $\mathbf{B} = \mathbf{I}_m$ and $\mathbf{C} = \mathbf{A}$
- ▶ If, on the other hand, m > n and rank(A) = n, then we can take B = A and C = I_n
- Example: Let $A = \begin{bmatrix} 2 & 1 & -2 & 5 \\ 1 & 0 & -3 & 2 \\ 3 & -1 & -13 & 5 \end{bmatrix} \quad rank(A) = 2$

We can write a full-rank factorization of *A* based on the proof of Lemma 12.3

$$\boldsymbol{A} = \begin{bmatrix} 2 & 1 \\ 1 & 0 \\ 3 & -1 \end{bmatrix} \begin{bmatrix} 1 & 0 & -3 & 2 \\ 0 & 1 & 4 & 1 \end{bmatrix} = \boldsymbol{B}\boldsymbol{C}$$

- Consider the matrix equation AXA = A, where A ∈ R^{m×n} is a given matrix and X ∈ R^{n×m} is a matrix we wish to determine.
 Observe that if A is a nonsingular square matrix, then the equation above has the unique solution X = A⁻¹
- ▶ Definition 12.1: Given A ∈ R^{m×n}, a matrix A[†] ∈ R^{n×m} is called a *pseudoinverse* of the matrix A if AA[†]A = A, and there exist matrices U ∈ R^{n×n} and V ∈ R^{m×m} such that A[†] = UA^T and A[†] = A^TV

- The requirement A[†] = UA^T = A^TV can be interpreted as follows. Each row of the pseudoinverse matrix A[†] of A is a linear combination of the rows of A^T, and each column of A[†] is a linear combination of the columns of A^T
- For the case which a matrix A ∈ R^{m×n} with m ≥ n and rank(A) = n, we can easily check that the following is a pseudoinverse of A:

$$\boldsymbol{A}^{\dagger} = (\boldsymbol{A}^T \boldsymbol{A})^{-1} \boldsymbol{A}^T$$

- Indeed, $A(A^T A)^{-1}A^T A = A$, and if we define $U = (A^T A)^{-1}$ and $V = A(A^T A)^{-1}(A^T A)^{-1}A^T$, then $A^{\dagger} = UA^T = A^T V$
- Note that we have A[†]A = I_n. For this reason, (A^TA)⁻¹A^T is often called the *left pseudoinverse* of A. This formula also appears in least-squares analysis (Sec. 12.1)

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For the case which a matrix A ∈ R^{m×n} with m ≤ n and rank(A) = m, we can easily check that the following is a pseudoinverse of A:

$$\boldsymbol{A}^{\dagger} = \boldsymbol{A}^T (\boldsymbol{A} \boldsymbol{A}^T)^{-1}$$

Note that we have AA[†] = I_m. For this reason, A^T(AA^T)⁻¹ is often called the *right pseudoinverse* of A. This formula also appears in the problem of minimizing ||x|| subject to Ax = b

- ▶ Theorem 12.4: Let A ∈ R^{m×n}. If a pseudoinverse A[†] of A exists, then it is unique.
- Our goal now is to show that the pseudoinverse matrix always exists. In fact, we show that the pseudoinverse of any given matrix *A* is given by the formula

$$oldsymbol{A}^{\dagger}=oldsymbol{C}^{\dagger}oldsymbol{B}^{\dagger}$$

where B^{\dagger} and C^{\dagger} are the pseudoinverse of the matrices B and C that form a full-rank factorization of A; that is, A = BC, where B and C are of full rank (Lemma 12.3)

Note that we already know how to compute B^{\dagger} and C^{\dagger} : $B^{\dagger} = (B^T B)^{-1} B^T$ $C^{\dagger} = C^T (C C^T)^{-1}$

• Theorem 12.5: Let a matrix $A \in \mathbb{R}^{m \times n}$ have a full-rank factorization A = BC, with rank(A) = rank(B) = rank(C) = r $B \in \mathbb{R}^{m \times r}$, $C \in \mathbb{R}^{r \times n}$, then

$$oldsymbol{A}^\dagger = oldsymbol{C}^\dagger oldsymbol{B}^\dagger$$

• Example:

(Does not necessarily hold if A = BC is not a full-rank factorization)

$$\boldsymbol{A} = \begin{bmatrix} 2 & 1 & -2 & 5 \\ 1 & 0 & -3 & 2 \\ 3 & -1 & -13 & 5 \end{bmatrix} = \begin{bmatrix} 2 & 1 \\ 1 & 0 \\ 3 & -1 \end{bmatrix} \begin{bmatrix} 1 & 0 & -3 & 2 \\ 0 & 1 & 4 & 1 \end{bmatrix} = \boldsymbol{B}\boldsymbol{C}$$
$$\boldsymbol{B}^{\dagger} = (\boldsymbol{B}^{T}\boldsymbol{B})^{-1}\boldsymbol{B}^{T} = \frac{1}{27} \begin{bmatrix} 5 & 2 & 5 \\ 16 & 1 & -11 \end{bmatrix}$$
$$\boldsymbol{C}^{\dagger} = \boldsymbol{C}^{T}(\boldsymbol{C}\boldsymbol{C}^{T})^{-1} = \frac{1}{76} \begin{bmatrix} 9 & 5 \\ 5 & 7 \\ -7 & 13 \\ 23 & 17 \end{bmatrix} \qquad \boldsymbol{A}^{\dagger} = \boldsymbol{C}^{\dagger}\boldsymbol{B}^{\dagger} = \frac{1}{2052} \begin{bmatrix} 123 & 23 & -10 \\ 137 & 17 & -52 \\ 173 & -1 & -178 \\ 387 & 63 & -72 \end{bmatrix}$$

• We can simplify the expression $A^{\dagger} = C^{\dagger}B^{\dagger} = C^{T}(CC^{T})^{-1}(B^{T}B)^{-1}B^{T}$ to $A^{\dagger} = C^{T}(B^{T}AA^{T})^{-1}B^{T}$

This is easily verified by substituting A = BC

▶ Theorem 12.6: Consider a system of linear equations Ax = b A ∈ R^{m×n}, rank(A) = r. The vector x* = A[†]b minimizes ||Ax - b||² on Rⁿ. Furthermore, among all vectors in Rⁿ that minimizes ||Ax - b||², the vector x* = A[†]b is the unique vector with minimal norm.

• The generalized inverse has the following useful properties

$$(\boldsymbol{A}^T)^{\dagger} = (\boldsymbol{A}^{\dagger})^T$$

$$(\boldsymbol{A}^{\dagger})^{\dagger} = \boldsymbol{A}$$

These two properties are similar to those that are satisfied by the usual matrix inverse. However, the property (A₁A₂)[†] = A[†]₂A[†]₁ does not hold in general.