

Math 5760/6890: Introduction to Mathematical Finance

Review: linear algebra and differential equation

Akil Narayan¹

¹Department of Mathematics, and Scientific Computing and Imaging (SCI) Institute
University of Utah

Fall 2024



We've discussed the basics of finance and investing – concepts of interest and present value.

A more advanced understanding of pricing and policies requires some math:

- **linear algebra**
- **differential equations**
- probability

These topics are prerequisites for this course!

Let $m, n \in \mathbb{N}$. ($m > n$, $m = n$, $m < n$ are allowed.)

We'll typically use lowercase boldface letters, e.g., \mathbf{v} , to denote *vectors*, elements of \mathbb{R}^n .

Vectors can be described by their components:

$$\mathbf{v} = \begin{pmatrix} v_1 \\ v_2 \\ \vdots \\ v_n \end{pmatrix} = \sum_{j=1}^n v_j \mathbf{e}_j \in \mathbb{R}^n, \quad \mathbf{e}_j = \begin{pmatrix} \vdots \\ 0 \\ 1 \\ 0 \\ \vdots \end{pmatrix}.$$

I.e., the components v_j are the *coordinates* of \mathbf{v} in an expansion of the canonical vectors $\{\mathbf{e}_j\}_{j=1}^n$.

Let $m, n \in \mathbb{N}$. ($m > n$, $m = n$, $m < n$ are allowed.)

We'll typically use lowercase boldface letters, e.g., \mathbf{v} , to denote *vectors*, elements of \mathbb{R}^n .

Vectors can be described by their components:

$$\mathbf{v} = \begin{pmatrix} v_1 \\ v_2 \\ \vdots \\ v_n \end{pmatrix} = \sum_{j=1}^n v_j \mathbf{e}_j \in \mathbb{R}^n, \quad \mathbf{e}_j = \begin{pmatrix} \vdots \\ 0 \\ 1 \\ 0 \\ \vdots \end{pmatrix}.$$

I.e., the components v_j are the *coordinates* of \mathbf{v} in an expansion of the canonical vectors $\{\mathbf{e}_j\}_{j=1}^n$.

We'll use uppercase boldface letters, e.g., \mathbf{A} , to denote *matrices*, elements of $\mathbb{R}^{m \times n}$ that are also described by their components:

$$\mathbf{A} = \begin{pmatrix} a_{11} & a_{12} & \cdots & a_{1n} \\ a_{21} & a_{22} & \cdots & a_{2n} \\ \vdots & & \ddots & \vdots \\ a_{m1} & a_{m2} & \cdots & a_{mn} \end{pmatrix} \in \mathbb{R}^{m \times n}.$$

Matrices are *linear maps* (functions) taking \mathbb{R}^n to \mathbb{R}^m .

It is sometimes useful to consider vectors as specializations of matrices:

- If $n = 1$ and $m > 1$, then $\mathbf{A} \in \mathbb{R}^{m \times 1}$ is a *column vector*
- If $m = 1$ and $n > 1$, then $\mathbf{A} \in \mathbb{R}^{1 \times n}$ is a *row vector*

When considering vectors as specializations of matrices, we will assume that vectors are column vectors, unless otherwise indicated.

Example (Portfolio parameterization)

Suppose we have some initial amount of money, $V(0)$, that we wish to invest.

Suppose there are $N \in \mathbb{N}$ securities, which are financial products of which we can purchase a quantity.

The price (per unit) of security i at time t is given by $S_i(t)$.

The number of units we purchase of security i is n_i (can be non-integer).

The *weight* of our portfolio for the i th security is $w_i = n_i S_i(0)/V(0)$, which is the relative amount of worth we invest in security i .

We represent all these things as vectors:

$$\mathbf{S}(t) = \begin{pmatrix} S_1(t) \\ \vdots \\ S_N(t) \end{pmatrix} \in \mathbb{R}^N, \quad \mathbf{n} = \begin{pmatrix} n_1 \\ \vdots \\ n_N \end{pmatrix} \in \mathbb{R}^N, \quad \mathbf{w} = \begin{pmatrix} w_1 \\ \vdots \\ w_N \end{pmatrix} \in \mathbb{R}^N.$$

The vector \mathbf{n} is the “trading strategy”, and \mathbf{w} is the (portfolio) “weight” vector.

The space of vectors \mathbb{R}^n has *Euclidean* structure. One source of this structure comes from the notion of inner products: With $\mathbf{v}, \mathbf{w} \in \mathbb{R}^n$, then the **inner product** of these vectors is

$$\langle \mathbf{v}, \mathbf{w} \rangle = \sum_{j=1}^n v_j w_j.$$

The inner product allows us to define lengths of vectors:

$$\|\mathbf{v}\| := \sqrt{\langle \mathbf{v}, \mathbf{v} \rangle} \geq 0,$$

with $\|\mathbf{v}\| = 0$ iff $\mathbf{v} = \mathbf{0}$.

The space of vectors \mathbb{R}^n has *Euclidean* structure. One source of this structure comes from the notion of inner products: With $\mathbf{v}, \mathbf{w} \in \mathbb{R}^n$, then the **inner product** of these vectors is

$$\langle \mathbf{v}, \mathbf{w} \rangle = \sum_{j=1}^n v_j w_j.$$

The inner product allows us to define lengths of vectors:

$$\|\mathbf{v}\| := \sqrt{\langle \mathbf{v}, \mathbf{v} \rangle} \geq 0,$$

with $\|\mathbf{v}\| = 0$ iff $\mathbf{v} = \mathbf{0}$.

From the definition, we observe that the inner product satisfies some key properties:

- *Symmetry*: $\langle \mathbf{v}, \mathbf{w} \rangle = \langle \mathbf{w}, \mathbf{v} \rangle$.
- *Bilinearity*: $\langle a\mathbf{u} + b\mathbf{v}, \mathbf{w} \rangle = a\langle \mathbf{u}, \mathbf{w} \rangle + b\langle \mathbf{v}, \mathbf{w} \rangle$ for any $a, b \in \mathbb{R}$.

A useful concept that inner products provide is a measure of angles between vectors:

$$\theta := \angle(\mathbf{v}, \mathbf{w}), \quad \cos \theta = \frac{\langle \mathbf{v}, \mathbf{w} \rangle}{\|\mathbf{v}\| \|\mathbf{w}\|}, \quad \mathbf{v}, \mathbf{w} \neq \mathbf{0}.$$

In particular this allows us to define *orthogonal* vectors: \mathbf{v} is orthogonal to \mathbf{w} if $\langle \mathbf{v}, \mathbf{w} \rangle = 0$.

A useful concept that inner products provide is a measure of angles between vectors:

$$\theta := \angle(\mathbf{v}, \mathbf{w}), \quad \cos \theta = \frac{\langle \mathbf{v}, \mathbf{w} \rangle}{\|\mathbf{v}\| \|\mathbf{w}\|}, \quad \mathbf{v}, \mathbf{w} \neq \mathbf{0}.$$

In particular this allows us to define *orthogonal* vectors: \mathbf{v} is orthogonal to \mathbf{w} if $\langle \mathbf{v}, \mathbf{w} \rangle = 0$.

Why should $\frac{\langle \mathbf{v}, \mathbf{w} \rangle}{\|\mathbf{v}\| \|\mathbf{w}\|}$ be a number between -1 and 1? Recall:

$$\left\langle \mathbf{v}, \frac{\mathbf{w}}{\|\mathbf{w}\|} \right\rangle = \text{“Amount” of } \mathbf{v} \text{ pointing in the direction of } \mathbf{w}.$$
$$\left\langle \mathbf{v}, \frac{\mathbf{w}}{\|\mathbf{w}\|} \right\rangle \mathbf{w} = \text{The projection of } \mathbf{v} \text{ onto } \mathbf{w}$$

A useful concept that inner products provide is a measure of angles between vectors:

$$\theta := \angle(\mathbf{v}, \mathbf{w}), \quad \cos \theta = \frac{\langle \mathbf{v}, \mathbf{w} \rangle}{\|\mathbf{v}\| \|\mathbf{w}\|}, \quad \mathbf{v}, \mathbf{w} \neq \mathbf{0}.$$

In particular this allows us to define *orthogonal* vectors: \mathbf{v} is orthogonal to \mathbf{w} if $\langle \mathbf{v}, \mathbf{w} \rangle = 0$.

Why should $\frac{\langle \mathbf{v}, \mathbf{w} \rangle}{\|\mathbf{v}\| \|\mathbf{w}\|}$ be a number between -1 and 1? Recall:

$$\left\langle \mathbf{v}, \frac{\mathbf{w}}{\|\mathbf{w}\|} \right\rangle = \text{“Amount” of } \mathbf{v} \text{ pointing in the direction of } \mathbf{w}.$$

$$\left\langle \mathbf{v}, \frac{\mathbf{w}}{\|\mathbf{w}\|} \right\rangle \mathbf{w} = \text{The projection of } \mathbf{v} \text{ onto } \mathbf{w}$$

If the first expression is the “amount” of \mathbf{v} pointing in a direction, then this “amount” shouldn’t be larger than $\|\mathbf{v}\|$:

$$\left| \left\langle \mathbf{v}, \frac{\mathbf{w}}{\|\mathbf{w}\|} \right\rangle \right| \leq \|\mathbf{v}\| \implies |\langle \mathbf{v}, \mathbf{w} \rangle| \leq \|\mathbf{v}\| \|\mathbf{w}\|$$

This is the **Cauchy-Schwarz** inequality. (Equality iff \mathbf{v} is a scalar multiple of \mathbf{w} .)

Because of Cauchy-Schwarz, the quantity $\frac{\langle \mathbf{v}, \mathbf{w} \rangle}{\|\mathbf{v}\| \|\mathbf{w}\|} \in [-1, 1]$, so that it can be the cosine of some angle.

Example

With a portfolio weight vector \mathbf{w} , the trading strategy \mathbf{n} , the per-unit security price $\mathbf{S}(t)$, and the initial capital $V(0)$, we have the following relations:

$$\langle \mathbf{w}, \mathbf{1} \rangle = \sum_{j=1}^N w_j = 1.$$

$$\langle \mathbf{n}, \mathbf{S}(0) \rangle = \sum_{j=1}^N n_j S_j(0) = V(0)$$

Example

With a portfolio weight vector \mathbf{w} , the trading strategy \mathbf{n} , the per-unit security price $\mathbf{S}(t)$, and the initial capital $V(0)$, we have the following relations:

$$\langle \mathbf{w}, \mathbf{1} \rangle = \sum_{j=1}^N w_j = 1.$$

$$\langle \mathbf{n}, \mathbf{S}(0) \rangle = \sum_{j=1}^N n_j S_j(0) = V(0)$$

There is no restriction on the values of the weights w_i : they can be negative or greater than 1.

- $w_i > 0$ corresponds to purchasing units, with the intention to sell later (a long position)
- $w_i < 0$ corresponds to borrowing units and selling them now, with the intention to buy them back later (“short selling”, a short position)

If there is no short selling, then $w_i \geq 0$, and hence $0 \leq w_i \leq 1$ for all i .

A core concept we'll need involves algebra on matrices, specifically *matrix multiplication*:

Given matrices $\mathbf{A} \in \mathbb{R}^{m \times \ell}$ and $\mathbf{B} \in \mathbb{R}^{\ell \times n}$, then the product \mathbf{AB} is given by,

$$\mathbf{AB} \in \mathbb{R}^{m \times n}, \quad (\mathbf{AB})_{j,k} = \sum_{q=1}^{\ell} A_{j,q} B_{q,k}$$

I.e., $(\mathbf{AB})_{j,k}$ is the inner product between the j th row of \mathbf{A} and the k th row of \mathbf{B} .

Matrix multiplication is defined for matrices of *conforming* sizes, i.e., when the inner dimensions match.

Matrix multiplication is in general *not* commutative.

A core concept we'll need involves algebra on matrices, specifically *matrix multiplication*:

Given matrices $\mathbf{A} \in \mathbb{R}^{m \times \ell}$ and $\mathbf{B} \in \mathbb{R}^{\ell \times n}$, then the product \mathbf{AB} is given by,

$$\mathbf{AB} \in \mathbb{R}^{m \times n}, \quad (\mathbf{AB})_{j,k} = \sum_{q=1}^{\ell} A_{j,q} B_{q,k}$$

i.e., $(\mathbf{AB})_{j,k}$ is the inner product between the j th row of \mathbf{A} and the k th row of \mathbf{B} .

Matrix multiplication is defined for matrices of *conforming* sizes, i.e., when the inner dimensions match.

Matrix multiplication is in general *not* commutative.

Given $\mathbf{A} \in \mathbb{R}^{m \times n}$, the *transpose* of \mathbf{A} is the matrix $\mathbf{A}^T \in \mathbb{R}^{n \times m}$, formed by reflecting the entries of \mathbf{A} across its main diagonal.

An inner product can be viewed as matrix multiplication:

$$\mathbf{v}^T \mathbf{w} = \langle \mathbf{v}, \mathbf{w} \rangle, \quad \mathbf{v}, \mathbf{w} \in \mathbb{R}^n.$$

(Recall that when interpreting vectors $\mathbf{v} \in \mathbb{R}^n$ as matrices, we consider them as column vectors $\mathbf{v} \in \mathbb{R}^{n \times 1}$).

An *outer product* is another matrix multiplication between vectors, but this time when the inner dimension is 1:

$$\mathbf{v} = (v_1, \dots, v_n)^T \in \mathbb{R}^n,$$

$$\mathbf{w} = (w_1, \dots, w_n)^T \in \mathbb{R}^n.$$

$$\mathbf{vw}^T = \left(\begin{array}{c|c|c} & & \\ \hline & w_1\mathbf{v} & w_2\mathbf{v} & \cdots & w_n\mathbf{v} \\ \hline & & & & \end{array} \right) = \left(\begin{array}{c|c|c} - & v_1\mathbf{w}^T & - \\ \hline - & v_2\mathbf{w}^T & - \\ & \vdots & \\ \hline - & v_n\mathbf{w}^T & - \end{array} \right) \in \mathbb{R}^{n \times n}$$

Let $\mathbf{v}_1, \dots, \mathbf{v}_k \in \mathbb{R}^n$ be any collection of vectors, and let $\mathbf{V} \in \mathbb{R}^{n \times k}$ be the matrix whose columns are these vectors:

$$\mathbf{V} = \left(\begin{array}{c|c|ccc|c} & & & & & \\ & & & & & \\ & & & & & \\ \mathbf{v}_1 & \mathbf{v}_2 & \cdots & \mathbf{v}_k & & \\ & & & & & \\ & & & & & \end{array} \right)$$

These vectors are **linearly dependent** if there exists a(ny) vector $\mathbf{c} \in \mathbb{R}^k$, $\mathbf{c} \neq \mathbf{0}$, such that,

$$\mathbf{V}\mathbf{c} = c_1\mathbf{v}_1 + \dots + c_k\mathbf{v}_k = \mathbf{0}.$$

Vectors that are *not* linearly dependent are **linearly independent**.

Vectors that are linearly dependent have a nontrivial linear relationship.

(If $\mathbf{0}$ is in the collection of vectors, the definition above implies they are linearly dependent.)

Let $\mathbf{v}_1, \dots, \mathbf{v}_k \in \mathbb{R}^n$ be any collection of vectors, and let $\mathbf{V} \in \mathbb{R}^{n \times k}$ be the matrix whose columns are these vectors:

$$\mathbf{V} = \left(\begin{array}{c|c|ccc|c} & & & & & \\ & & & & & \\ & & & & & \\ \mathbf{v}_1 & \mathbf{v}_2 & \cdots & & & \mathbf{v}_k \\ & & & & & \\ & & & & & \\ & & & & & \end{array} \right)$$

The **span** of these vectors is the collection of all linear combinations of these vectors:

$$\text{span}\{\mathbf{v}_1, \dots, \mathbf{v}_k\} := \{\mathbf{V}\mathbf{c} \mid \mathbf{c} \in \mathbb{R}^k\}.$$

The span of vectors is a *linear/vector subspace*: it is a collection of vectors closed under addition and scalar multiplication.

Let $\mathbf{v}_1, \dots, \mathbf{v}_k \in \mathbb{R}^n$ be any collection of vectors, and let $\mathbf{V} \in \mathbb{R}^{n \times k}$ be the matrix whose columns are these vectors:

$$\mathbf{V} = \left(\begin{array}{c|c|ccc|c} & & & & & \\ & & & & & \\ & & & & & \\ \mathbf{v}_1 & \mathbf{v}_2 & \cdots & \mathbf{v}_k & & \\ & & & & & \end{array} \right)$$

Let S be some given vector subspace.

The vectors form a **basis** for S if the span of these vectors is S and they are linearly independent.

Let $\mathbf{v}_1, \dots, \mathbf{v}_k \in \mathbb{R}^n$ be any collection of vectors, and let $\mathbf{V} \in \mathbb{R}^{n \times k}$ be the matrix whose columns are these vectors:

$$\mathbf{V} = \left(\begin{array}{c|c|ccc|c} & & & & & \\ & & & & & \\ & & & & & \\ \mathbf{v}_1 & \mathbf{v}_2 & \cdots & \mathbf{v}_k & & \\ & & & & & \end{array} \right)$$

Let S be some given vector subspace.

The vectors form a **basis** for S if the span of these vectors is S and they are linearly independent.

In math: these vectors are a basis for S if

$$\forall \mathbf{w} \in S, \exists! \mathbf{c} \in \mathbb{R}^k \text{ such that } \mathbf{V}\mathbf{c} = \mathbf{w}.$$

(If \mathbf{c} did not exist, the vectors wouldn't span S . If \mathbf{c} weren't unique, then there would exist a nontrivial solution to $\mathbf{V}\mathbf{d} = \mathbf{0}$.)

Let $\mathbf{v}_1, \dots, \mathbf{v}_k \in \mathbb{R}^n$ be any collection of vectors, and let $\mathbf{V} \in \mathbb{R}^{n \times k}$ be the matrix whose columns are these vectors:

$$\mathbf{V} = \left(\begin{array}{c|c|ccc|c} & & & & & \\ & & & & & \\ & & & & & \\ \mathbf{v}_1 & \mathbf{v}_2 & \cdots & \mathbf{v}_k & & \\ & & & & & \end{array} \right)$$

Let S be some given vector subspace.

The vectors form a **basis** for S if the span of these vectors is S and they are linearly independent.

In math: these vectors are a basis for S if

$$\forall \mathbf{w} \in S, \exists! \mathbf{c} \in \mathbb{R}^k \text{ such that } \mathbf{V}\mathbf{c} = \mathbf{w}.$$

(If \mathbf{c} did not exist, the vectors wouldn't span S . If \mathbf{c} weren't unique, then there would exist a nontrivial solution to $\mathbf{V}\mathbf{d} = \mathbf{0}$.)

A basis for S is not unique, but the size of a basis for S is unique.

This unique size of a basis for S is its **dimension**, $\dim S$.

If S contains m -dimensional vectors, then $\dim S \leq m$.

One particularly important application of linear algebra is as the theoretical and practical underpinning for solving linear equations for an unknown vector $\mathbf{x} \in \mathbb{R}^n$:

$$\mathbf{Ax} = \mathbf{b}, \quad \mathbf{A} = \begin{pmatrix} | & | & \cdots & | \\ \mathbf{a}_1 & \mathbf{a}_2 & \cdots & \mathbf{a}_n \\ | & | & \cdots & | \end{pmatrix} \in \mathbb{R}^{m \times n}, \quad \mathbf{b} \in \mathbb{R}^m.$$

To characterize solutions to such linear equations, consider the *range* or “column space” of \mathbf{A} , which is a subspace:

$$\text{range}(\mathbf{A}) := \text{span}\{\mathbf{a}_1, \dots, \mathbf{a}_n\} \implies n \geq \dim \text{range}(\mathbf{A})$$

One particularly important application of linear algebra is as the theoretical and practical underpinning for solving linear equations for an unknown vector $\mathbf{x} \in \mathbb{R}^n$:

$$\mathbf{Ax} = \mathbf{b}, \quad \mathbf{A} = \begin{pmatrix} | & | & \cdots & | \\ \mathbf{a}_1 & \mathbf{a}_2 & \cdots & \mathbf{a}_n \\ | & | & & | \end{pmatrix} \in \mathbb{R}^{m \times n}, \quad \mathbf{b} \in \mathbb{R}^m.$$

To characterize solutions to such linear equations, consider the *range* or “column space” of \mathbf{A} , which is a subspace:

$$\text{range}(\mathbf{A}) := \text{span}\{\mathbf{a}_1, \dots, \mathbf{a}_n\} \implies n \geq \dim \text{range}(\mathbf{A})$$

We can make very strong characterizations about solutions to linear systems:

1. If $\mathbf{b} \notin \text{range}(\mathbf{A})$, then there is no solution \mathbf{x} .
2. If $\mathbf{b} \in \text{range}(\mathbf{A})$ and $n > \dim \text{range}(\mathbf{A})$ then there are infinitely many solutions \mathbf{x} , and the collection of these solutions form an *affine space*¹ of dimension $(n - \dim \text{range}(\mathbf{A}))$.
3. If $\mathbf{b} \in \text{range}(\mathbf{A})$ and $n = \dim \text{range}(\mathbf{A})$, then there exists exactly one solution \mathbf{x} .

¹An affine space is a subspace shifted by a fixed vector.

One particularly important application of linear algebra is as the theoretical and practical underpinning for solving linear equations for an unknown vector $\mathbf{x} \in \mathbb{R}^n$:

$$\mathbf{Ax} = \mathbf{b}, \quad \mathbf{A} = \begin{pmatrix} | & | & \cdots & | \\ \mathbf{a}_1 & \mathbf{a}_2 & \cdots & \mathbf{a}_n \\ | & | & & | \end{pmatrix} \in \mathbb{R}^{m \times n}, \quad \mathbf{b} \in \mathbb{R}^m.$$

To characterize solutions to such linear equations, consider the *range* or “column space” of \mathbf{A} , which is a subspace:

$$\text{range}(\mathbf{A}) := \text{span}\{\mathbf{a}_1, \dots, \mathbf{a}_n\} \implies n \geq \dim \text{range}(\mathbf{A})$$

We can make very strong characterizations about solutions to linear systems:

1. If $\mathbf{b} \notin \text{range}(\mathbf{A})$, then there is no solution \mathbf{x} .
2. If $\mathbf{b} \in \text{range}(\mathbf{A})$ and $n > \dim \text{range}(\mathbf{A})$ then there are infinitely many solutions \mathbf{x} , and the collection of these solutions form an *affine space*¹ of dimension $(n - \dim \text{range}(\mathbf{A}))$.
3. If $\mathbf{b} \in \text{range}(\mathbf{A})$ and $n = \dim \text{range}(\mathbf{A})$, then there exists exactly one solution \mathbf{x} .

NB: Situations 1 and 2 can happen for any relationship between n and m . Situation 3 can happen only if $m \geq n$. The canonical algorithm to compute solutions to linear equations is Gaussian elimination.

¹An affine space is a subspace shifted by a fixed vector.

Example

Recall that portfolio weights satisfy,

$$\langle \mathbf{w}, \mathbf{1} \rangle = 1.$$

This is equivalent to:

$$\mathbf{A}\mathbf{w} = \mathbf{b}, \quad \mathbf{A} = \begin{pmatrix} 1 & 1 & \cdots & 1 \end{pmatrix} \in \mathbb{R}^{1 \times N}, \quad \mathbf{b} = 1 \in \mathbb{R}^1.$$

Example

Recall that portfolio weights satisfy,

$$\langle \mathbf{w}, \mathbf{1} \rangle = 1.$$

This is equivalent to:

$$\mathbf{A}\mathbf{w} = \mathbf{b}, \quad \mathbf{A} = \begin{pmatrix} 1 & 1 & \cdots & 1 \end{pmatrix} \in \mathbb{R}^{1 \times N}, \quad \mathbf{b} = 1 \in \mathbb{R}^1.$$

In this case, the dimension of the range is $\dim \text{range}(\mathbf{A}) = 1$ (and clearly $\mathbf{b} \in \text{range}(\mathbf{A})$).

Hence, there are infinitely many valid portfolio weight vectors \mathbf{w} , and they form an affine space of dimension $N - \dim \text{range}(\mathbf{A}) = N - 1$.

When $m = n$, consider the “square” linear system,

$$\mathbf{A}\mathbf{x} = \mathbf{b},$$

\mathbf{A}, \mathbf{b} given

There are some equivalent statements about a unique solution:

- There is a unique solution \mathbf{x} .
- The **rank** of \mathbf{A} , that is $\dim \text{range}(\mathbf{A})$, has maximal value n .
- The *determinant* of \mathbf{A} does not vanish: $\det \mathbf{A} \neq 0$.
- The matrix \mathbf{A} has an *inverse* \mathbf{A}^{-1} , satisfying $\mathbf{A}\mathbf{A}^{-1} = \mathbf{A}^{-1}\mathbf{A} = \mathbf{I}$.

When $m = n$, consider the “square” linear system,

$$\mathbf{A}\mathbf{x} = \mathbf{b},$$

\mathbf{A}, \mathbf{b} given

There are some equivalent statements about a unique solution:

- There is a unique solution \mathbf{x} .
- The **rank** of \mathbf{A} , that is $\dim \text{range}(\mathbf{A})$, has maximal value n .
- The *determinant* of \mathbf{A} does not vanish: $\det \mathbf{A} \neq 0$.
- The matrix \mathbf{A} has an *inverse* \mathbf{A}^{-1} , satisfying $\mathbf{A}\mathbf{A}^{-1} = \mathbf{A}^{-1}\mathbf{A} = \mathbf{I}$.

When any (hence all) of the above statements is true, then

$$\mathbf{x} = \mathbf{A}^{-1}\mathbf{b}$$

is the unique solution.

Orthogonal matrices

Matrix inverses are generally “hard” to compute (analytically or numerically).

There is one class of matrices for which matrix inversion is rather simple:

A matrix $\mathbf{A} \in \mathbb{R}^{n \times n}$ is **orthogonal** if its columns are (pairwise) orthogonal and unit norm:

$$\mathbf{A} = \left(\begin{array}{c|c|c|c} | & | & \cdots & | \\ \mathbf{a}_1 & \mathbf{a}_2 & & \mathbf{a}_n \\ | & | & & | \end{array} \right), \quad \langle \mathbf{a}_j, \mathbf{a}_k \rangle = \delta_{j,k} := \begin{cases} 1, & j = k \\ 0, & j \neq k \end{cases}$$

Orthogonal matrices

Matrix inverses are generally “hard” to compute (analytically or numerically).

There is one class of matrices for which matrix inversion is rather simple:

A matrix $\mathbf{A} \in \mathbb{R}^{n \times n}$ is **orthogonal** if its columns are (pairwise) orthogonal and unit norm:

$$\mathbf{A} = \left(\begin{array}{c|c|c|c} | & | & \cdots & | \\ \mathbf{a}_1 & \mathbf{a}_2 & & \mathbf{a}_n \\ | & | & & | \end{array} \right), \quad \langle \mathbf{a}_j, \mathbf{a}_k \rangle = \delta_{j,k} := \begin{cases} 1, & j = k \\ 0, & j \neq k \end{cases}$$

A straightforward computation using matrix multiplication reveals:

$$\mathbf{A} \text{ orthogonal} \implies \mathbf{A}^T \mathbf{A} = \mathbf{I} \implies \mathbf{A}^{-1} = \mathbf{A}^T.$$

Hence, orthogonality is a particularly useful practical property. (And \mathbf{A} orthogonal implies $\mathbf{A}^{-1} = \mathbf{A}^T$ is also orthogonal.)

Orthogonal matrices

Matrix inverses are generally “hard” to compute (analytically or numerically).

There is one class of matrices for which matrix inversion is rather simple:

A matrix $\mathbf{A} \in \mathbb{R}^{n \times n}$ is **orthogonal** if its columns are (pairwise) orthogonal and unit norm:

$$\mathbf{A} = \left(\begin{array}{c|c|c|c} | & | & \cdots & | \\ \mathbf{a}_1 & \mathbf{a}_2 & & \mathbf{a}_n \\ | & | & & | \end{array} \right), \quad \langle \mathbf{a}_j, \mathbf{a}_k \rangle = \delta_{j,k} := \begin{cases} 1, & j = k \\ 0, & j \neq k \end{cases}$$

A straightforward computation using matrix multiplication reveals:

$$\mathbf{A} \text{ orthogonal} \implies \mathbf{A}^T \mathbf{A} = \mathbf{I} \implies \mathbf{A}^{-1} = \mathbf{A}^T.$$

Hence, orthogonality is a particularly useful practical property. (And \mathbf{A} orthogonal implies $\mathbf{A}^{-1} = \mathbf{A}^T$ is also orthogonal.)

Another useful property of orthogonal matrices: they correspond to *isometric* maps.

In particular, if \mathbf{A} is orthogonal:

$$\langle \mathbf{A}\mathbf{v}, \mathbf{A}\mathbf{w} \rangle = \mathbf{v}^T \mathbf{A}^T \mathbf{A} \mathbf{w} = \mathbf{v}^T \mathbf{I} \mathbf{w} = \mathbf{v}^T \mathbf{w} = \langle \mathbf{v}, \mathbf{w} \rangle.$$

I.e., the transformation $\mathbf{v} \mapsto \mathbf{A}\mathbf{v}$ preserves angles and lengths.

Orthogonal matrices are simple rotations and/or reflections.

For square matrices $\mathbf{A} \in \mathbb{R}^{n \times n}$, an important concept is that of the *spectrum* of \mathbf{A} .

If there exists any (possibly complex-valued) scaled λ , and any non-zero vector \mathbf{v} (possibly complex-valued) such that,

$$\mathbf{A}\mathbf{v} = \lambda\mathbf{v},$$

then

- λ is called an *eigenvalue* of \mathbf{A}
- \mathbf{v} is called an *eigenvector* of \mathbf{A} .

For square matrices $\mathbf{A} \in \mathbb{R}^{n \times n}$, an important concept is that of the *spectrum* of \mathbf{A} .

If there exists any (possibly complex-valued) scalar λ , and any non-zero vector \mathbf{v} (possibly complex-valued) such that,

$$\mathbf{A}\mathbf{v} = \lambda\mathbf{v},$$

then

- λ is called an *eigenvalue* of \mathbf{A}
- \mathbf{v} is called an *eigenvector* of \mathbf{A} .

Fixing λ , note that the condition for being an eigenvector is invariant under addition of vectors and scalar multiplication: The set of eigenvectors associated to an eigenvalue λ is a subspace.

For square matrices $\mathbf{A} \in \mathbb{R}^{n \times n}$, an important concept is that of the *spectrum* of \mathbf{A} .

If there exists any (possibly complex-valued) scaled λ , and any non-zero vector \mathbf{v} (possibly complex-valued) such that,

$$\mathbf{A}\mathbf{v} = \lambda\mathbf{v},$$

then

- λ is called an *eigenvalue* of \mathbf{A}
- \mathbf{v} is called an *eigenvector* of \mathbf{A} .

Fixing λ , note that the condition for being an eigenvector is invariant under addition of vectors and scalar multiplication: The set of eigenvectors associated to an eigenvalue λ is a subspace.

Eigenvalues λ satisfy the characteristic equation:

$$\det(\mathbf{A} - \lambda\mathbf{I}) = 0,$$

so that eigenvalues are roots of a degree- n polynomial.

Matrix diagonalization

D06-S19(a)

Every $n \times n$ matrix has exactly n eigenvalues (possibly repeated according to roots of the characteristic equation).

Every $n \times n$ matrix has exactly n eigenvalues (possibly repeated according to roots of the characteristic equation).

For each eigenvalue (counting multiplicity), there *may* be an eigenvector that is linearly independent from all others.

Matrices for which each eigenvalue has a corresponding linearly independent eigenvector are called *diagonalizable*.

If \mathbf{A} is diagonalizable, then the following decomposition holds,

$$\mathbf{A} = \mathbf{V}\mathbf{\Lambda}\mathbf{V}^{-1}, \quad \mathbf{\Lambda} = \text{diag}(\lambda_1, \dots, \lambda_n), \quad \mathbf{V} = \left(\begin{array}{c|c|c|c} | & | & \cdots & | \\ \mathbf{v}_1 & \mathbf{v}_2 & & \mathbf{v}_n \\ | & | & & | \end{array} \right),$$

where $(\lambda_j, \mathbf{v}_j)$ are eigenvalue-eigenvector pairs for $j = 1, \dots, n$.

A particularly nice property about diagonalizable matrices is that the eigenvectors span \mathbb{R}^n (possibly using complex scalar multiplication).

The upshot: if \mathbf{A} is diagonalizable, then there is a linear transformation (defined by \mathbf{V}) such that multiplication by \mathbf{A} corresponds to a simple diagonal scaling:

$$\mathbf{w} = \mathbf{A}\mathbf{x} \quad \xrightarrow{\mathbf{y}=\mathbf{V}^{-1}\mathbf{x}, \mathbf{z}=\mathbf{V}^{-1}\mathbf{w}} \quad \mathbf{z} = \mathbf{\Lambda}\mathbf{y}.$$

Hence diagonalizations are *very* useful.

Orthogonal diagonalization

D06-S20(a)

Diagonalizable matrices are diagonal under some transformation defined by \mathbf{V}^{-1} . But \mathbf{V}^{-1} can be painful to compute.

Some matrices are *orthogonally* diagonalizable, meaning that \mathbf{V} is an orthogonal matrix, and hence \mathbf{V}^{-1} is “easy” to compute.

Diagonalizable matrices are diagonal under some transformation defined by \mathbf{V}^{-1} . But \mathbf{V}^{-1} can be painful to compute.

Some matrices are *orthogonally* diagonalizable, meaning that \mathbf{V} is an orthogonal matrix, and hence \mathbf{V}^{-1} is “easy” to compute.

One of the major results of linear algebra is the following identification of one class of orthogonal matrices:

Theorem (Spectral theorem for symmetric matrices)

Assume $\mathbf{A} \in \mathbb{R}^{n \times n}$ satisfies $\mathbf{A} = \mathbf{A}^T$.
(Such matrices are called symmetric.)

Then:

- All eigenvalues of \mathbf{A} are real-valued.
- \mathbf{A} is orthogonally diagonalizable.
(The eigenvectors can be chosen as orthogonal vectors.)

Diagonalizable matrices are diagonal under some transformation defined by V^{-1} . But V^{-1} can be painful to compute.

Some matrices are *orthogonally* diagonalizable, meaning that V is an orthogonal matrix, and hence V^{-1} is “easy” to compute.

One of the major results of linear algebra is the following identification of one class of orthogonal matrices:

Theorem (Spectral theorem for symmetric matrices)

Assume $A \in \mathbb{R}^{n \times n}$ satisfies $A = A^T$.
(Such matrices are called symmetric.)

Then:

- All eigenvalues of A are real-valued.
- A is orthogonally diagonalizable.
(The eigenvectors can be chosen as orthogonal vectors.)

I.e.,

$$A = A^T \implies A = V\Lambda V^T.$$

Since symmetric matrices have real-valued eigenvalues, then one can make sensible definitions about where the eigenvalues lie on \mathbb{R} .

In particular, the following are well-defined for λ_j the eigenvalues of an $n \times n$ symmetric matrix:

- $\lambda_{\min} = \min_{j=1, \dots, n} \lambda_j$
- $\lambda_{\max} = \max_{j=1, \dots, n} \lambda_j$

Equality above occurs iff \boldsymbol{x} is a multiple of the eigenvalue corresponding to the minimum/maximum eigenvalue of \boldsymbol{A} .

Since symmetric matrices have real-valued eigenvalues, then one can make sensible definitions about where the eigenvalues lie on \mathbb{R} .

In particular, the following are well-defined for λ_j the eigenvalues of an $n \times n$ symmetric matrix:

- $\lambda_{\min} = \min_{j=1, \dots, n} \lambda_j$
- $\lambda_{\max} = \max_{j=1, \dots, n} \lambda_j$

Equality above occurs iff \mathbf{x} is a multiple of the eigenvalue corresponding to the minimum/maximum eigenvalue of \mathbf{A} .

The spectral theorem also implies the following extremely useful inequality: if \mathbf{A} is symmetric, then,

$$\lambda_{\min} \|\mathbf{x}\|^2 \leq \mathbf{x}^T \mathbf{A} \mathbf{x} \leq \lambda_{\max} \|\mathbf{x}\|^2.$$

(The function $f(\mathbf{x}) = \mathbf{x}^T \mathbf{A} \mathbf{x}$ is an example of a *quadratic form*.)

Two final definitions are sub-classes of symmetric matrices:

- If \mathbf{A} is symmetric and all its eigenvalues are strictly positive, then \mathbf{A} is (symmetric) **positive definite**.
- If \mathbf{A} is symmetric and all its eigenvalues are non-negative, then \mathbf{A} is (symmetric) **positive semidefinite**.

One can equivalently define these matrix classes through their quadratic forms.

In particular, \mathbf{A} is symmetric positive semi-definite iff $\mathbf{x}^T \mathbf{A} \mathbf{x} \geq 0$ for all $\mathbf{x} \in \mathbb{R}^n$.

Differential equations govern how quantities change in time.

One class of general ordinary differential equations (DE) governing the unknown function $y(t)$ where t is a scalar (i.e., time) is,

$$F(t, y, y', y'', y''', \dots) = 0, \quad y(0) = y_0, \quad y'(0) = y'_0 \quad \dots$$

This is an **initial value problem**. The maximum derivative appearing in F is called the *order* of the equation.

Differential equations govern how quantities change in time.

One class of general ordinary differential equations (DE) governing the unknown function $y(t)$ where t is a scalar (i.e., time) is,

$$F(t, y, y', y'', y''', \dots) = 0, \quad y(0) = y_0, \quad y'(0) = y'_0 \quad \dots$$

This is an **initial value problem**. The maximum derivative appearing in F is called the *order* of the equation.

Understanding the theory (solvability, well-posedness) of these problems is generally quite difficult, but *linear* equations are quite flexible for modeling and are rather well-understood.

Linear DE's are those where y , y' , etc., collectively appear in F in a *linear* fashion.

The initial value problem,

$$y'(t) = 3y,$$

$$y(0) = 4,$$

is, with some experience, rather transparent to solve:

$$y(t) = 4e^{3t}.$$

The initial value problem,

$$y'(t) = 3y, \qquad y(0) = 4,$$

is, with some experience, rather transparent to solve:

$$y(t) = 4e^{3t}.$$

The general solution of a first-order “constant-coefficient” linear equation can be determined in a similar fashion:

$$y' + \lambda y = f(t), \qquad y(0) = y_0,$$

has the unique solution

$$y(t) = y_0 e^{-\lambda t} + e^{-\lambda t} \int_0^t f(s) e^{\lambda s} ds$$

Example

Consider an annuity (say a loan) with interest continuously compounded at (annual) rate r .

Suppose we also consider a repayment cycle that is *continuous* instead of periodic. I.e., we pay money at a (continuous) rate of P dollars (per year).

Example

Consider an annuity (say a loan) with interest continuously compounded at (annual) rate r .

Suppose we also consider a repayment cycle that is *continuous* instead of periodic. I.e., we pay money at a (continuous) rate of P dollars (per year).

The equation modeling the time- t present value $V(t)$ of this annuity is given by,

$$V'(t) = rV(t) - P, \quad V(0) = V_0,$$

where V_0 is the loan principal (annuity present value at time 0).

Example

Consider an annuity (say a loan) with interest continuously compounded at (annual) rate r .

Suppose we also consider a repayment cycle that is *continuous* instead of periodic. I.e., we pay money at a (continuous) rate of P dollars (per year).

The equation modeling the time- t present value $V(t)$ of this annuity is given by,

$$V'(t) = rV(t) - P, \quad V(0) = V_0,$$

where V_0 is the loan principal (annuity present value at time 0).

The solution to this equation is given by,

$$\begin{aligned} V(t) &= V_0 e^{rt} - e^{rt} \int_0^t P e^{-rs} ds \\ &= V_0 e^{rt} + \frac{P}{r} [1 - e^{rt}] \end{aligned}$$

Example

Consider an annuity (say a loan) with interest continuously compounded at (annual) rate r .

Suppose we also consider a repayment cycle that is *continuous* instead of periodic. I.e., we pay money at a (continuous) rate of P dollars (per year).

The equation modeling the time- t present value $V(t)$ of this annuity is given by,

$$V'(t) = rV(t) - P, \quad V(0) = V_0,$$

where V_0 is the loan principal (annuity present value at time 0).

The solution to this equation is given by,

$$\begin{aligned} V(t) &= V_0 e^{rt} - e^{rt} \int_0^t P e^{-rs} ds \\ &= V_0 e^{rt} + \frac{P}{r} [1 - e^{rt}] \end{aligned}$$

Note in particular that this implies $V_0 < P/r$ is required in order for the loan to eventually be repaid.

The rather interesting part of this comes with *systems* of linear constant-coefficient differential equations:

$$\mathbf{y}'(t) = \mathbf{A}\mathbf{y},$$

$$\mathbf{y}(0) = \mathbf{y}_0,$$

where $\mathbf{y}(t) = (y_1(t), y_2(t), \dots, y_n(t))^T$.

The rather interesting part of this comes with *systems* of linear constant-coefficient differential equations:

$$\mathbf{y}'(t) = \mathbf{A}\mathbf{y}, \quad \mathbf{y}(0) = \mathbf{y}_0,$$

where $\mathbf{y}(t) = (y_1(t), y_2(t), \dots, y_n(t))^T$.

If \mathbf{A} is diagonalizable, $\mathbf{A} = \mathbf{V}\mathbf{\Lambda}\mathbf{V}^{-1}$, then this system is the same as

$$\mathbf{z}'(t) = \mathbf{\Lambda}\mathbf{z}, \quad \mathbf{z}(0) = \mathbf{V}^{-1}\mathbf{y}_0,$$

where $\mathbf{z}(t) := \mathbf{V}^{-1}\mathbf{y}(t)$, and easily solvable:

$$\mathbf{z}(t) = e^{\mathbf{\Lambda}t}\mathbf{V}^{-1}\mathbf{y}_0 \quad \implies \quad \mathbf{y}(t) = \mathbf{V}e^{\mathbf{\Lambda}t}\mathbf{V}^{-1}\mathbf{y}_0.$$

Above, $e^{\mathbf{\Lambda}t} = (e^{\lambda_1 t}, e^{\lambda_2 t}, \dots, e^{\lambda_n t})^T$.

(If \mathbf{A} is orthogonally diagonalizable, this is computationally even easier.)

