MATH 20700 (Honors Analysis in \mathbb{R}^n I) Notes

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Part I Linear Algebra

Basic Notions

1.1 Notes

9/27: • Vector space: Basically, a set for which you have an addition and multiplication.

- \mathbb{F}^d is used for \mathbb{R}^d or \mathbb{C}^d in Treil (2017).
- \mathbb{P}_n is the vector space of polynomials up to degree n.
- C([0,1]) is the set of continuous functions defined on [0,1], an infinite-dimensional vector space.
- Generating set: A subset of a vector space, all linear combinations of which generate the vector space. Also known as spanning set.
 - Any element of VS is a linear comb. of elements of the generating set.
- Linearly independent (list): A list of vectors $\mathbf{v}_1, \dots, \mathbf{v}_k \in V$ such that $\sum_{i=1}^k \alpha_i \mathbf{v}_i = 0$ implies $\alpha_i = 0$ for all i.
- Base: A generating set consisting of linearly independent vectors.
- Any element of a VS can be written as a unique linear combination of the vectors in a base.
 - If $\mathbf{x} = \sum_{i=1}^k \alpha_i \mathbf{v}_i = \sum_{i=1}^k \beta_i \mathbf{v}_i$, then $\alpha_i = \beta_i$ for all i.
- Linear transformation: A function $T: X \to Y$, where X, Y are VSs, such that

$$T(\alpha \mathbf{x} + \beta \mathbf{y}) = \alpha T \mathbf{x} + \beta T \mathbf{y}$$

for all $\mathbf{x} \in X$, $\mathbf{y} \in Y$.

- Examples of linear transformations:
 - Consider \mathbb{P}_n . Let $Tp_n = p'_n$. This T is linear.
 - Rotation in \mathbb{R}^d .
 - \blacksquare Think graphically about two vectors \mathbf{x}, \mathbf{y} .
 - Rotating and summing them is the same as summing and rotating. Same for scaling.
 - Thus, rotation is actually linear!
 - Reflection as well.
- Consider $T: \mathbb{R} \to \mathbb{R}$.
 - Any linear map on the line is a line.
 - We must have $Tx = \alpha x$: $Tx = T(1x) = xT(1) = x\alpha$.

- Consider $T: \mathbb{R}^n \to \mathbb{R}^m$ linear.
 - Any linear map between \mathbb{R}^n and \mathbb{R}^m is linear.
 - Thus, $T(\mathbf{x}) = A\mathbf{x}$ for all $\mathbf{x} \in \mathbb{R}^n$, where A is an $m \times n$ matrix.
- To find A, do the same calculation as for $Tx = \alpha x$ but more carefully:
 - Let $\{\mathbf{e}_1, \dots, \mathbf{e}_n\}$ be a basis.
 - So $\mathbf{x} = \sum_{i=1}^{n} \alpha_i \mathbf{e}_i$.
 - Thus, $T\mathbf{x} = \sum_{i=1}^{n} \alpha_i T(\mathbf{e}_i)$.
 - Each $T(\mathbf{e}_i)$ is part of the matrix that we multiply by the column vector representing \mathbf{x} .
- Multiplication of matrices is equivalent to composition of linear maps.
- Consider $T_1: \mathbb{R}^n \to \mathbb{R}^m$ and $T_2: \mathbb{R}^m \to \mathbb{F}^r$.
 - $T_2 \circ T_1$ is equivalent to BA, if A represents T_1 and B represents T_2 . In other words, $(T_2 \circ T_1)(\mathbf{x}) = BA\mathbf{x}$ for all \mathbf{x} .
- Recall that if $A = (\alpha_{ij})$ and $B = (\beta_{ij})$, then $(BA)_{ij} = (\sum \beta_{ik} \alpha_{kj})$.
- Properties of multiplication:

$$(AB)C = A(BC)$$
$$A(B+C) = AB + AC$$
$$(A+B)C = AC + BC$$

- However, it is not true in general that AB = BA.
- Trace (of an $n \times n$ matrix A): The sum of the diagonal entries of A. Denoted by tr(A). Given by

$$\operatorname{tr}(A) = \sum \alpha_{ii}$$

- It is true that tr(AB) = tr(BA).
 - Indeed, on the diagonals, multiplication is commutative; it's the other terms that mess you up in general.
- Invertibility of matrices.
 - In general, matrices are not invertible: Not every system of equations is solveable; Ax = b does not always have a solution $x = A^{-1}b$.
- C is the inverse from the left: CA = I. B is the inverse from the right: AB = I. A matrix can have a left and a right inverse and still not be invertible. A matrix is invertible iff C = B.
- Any time we write "inverse," we do so under the assumption that it exists.
- $(AB)^{-1} = B^{-1}A^{-1}$ easy proof by multiplication.
- If $A = (a_{ij}), A^T = (a_{ji}).$
 - $(A^{-1})^T = (A^T)^{-1}.$
 - $(AB)^T = B^T A^T.$
- Let X, Y VS.
 - $-X \cong Y^{[1]}$ if there exists a linear $T: X \to Y$ that is one-to-one and onto.
 - Check: A(basis of X) = basis of Y. Prove by definition and expression of elements as linear combinations.
- Subspace: A subset of a vector space which happens to be a vector space, itself.

 $^{^1}$ "X is isomorphic to Y."

1.2 Chapter 1: Basic Notions

From Treil (2017).

10/24:

- Coordinates (of $\mathbf{v} \in V$ wrt. a basis $\mathbf{v}_1, \dots, \mathbf{v}_n$ of V): The unique scalars $\alpha_1, \dots, \alpha_n \in \mathbb{F}$ such that $\mathbf{v} = \alpha_1 \mathbf{v}_1 + \dots + \alpha_n \mathbf{v}_n$.
- Spanning system: A list of vectors that spans V. Also known as generating system, complete system.
- Trivial (linear combination): A linear combination $\alpha_1 \mathbf{v}_1 + \cdots + \alpha_n \mathbf{v}_n$ of vectors such that $\alpha_k = 0$ for each $k = 1, \dots, n$.
- Transformation: A function $T: X \to Y$. Also known as transform, mapping, map, operator, and function.
- The matrix of a linear transformation T is often denoted by [T].
- To compute the reflection of vectors over an arbitrary line through the origin in \mathbb{R}^2 , represent the overall transformation as a composition of rotating the line to be the x-axis, reflecting over the x-axis, and rotating back.
- Theorem 1.5.1: If A is an $m \times n$ matrix and B is an $n \times m$ matrix, then

$$\operatorname{tr}(AB) = \operatorname{tr}(BA)$$

- Theorem 1.6.1: If a linear transformation is invertible, then its left and right inverses are unique and coincide.
- The column $(1,1)^T$ is left-invertible, with one possible left inverse being (1/2,1/2).
 - Note that it is not right invertible since its left inverses are not unique (see Theorem 1.6.1).
- An invertible matrix must be square.
- **Isomorphic** (vector spaces): Two vectors spaces V, W such that there exists an isomorphism $A: V \to W$. Denoted by $V \cong W$.
- Theorem 1.6.8: $A: X \to Y$ is invertible if and only if for any right side $\mathbf{b} \in Y$, the equation

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

has a unique solution $\mathbf{x} \in X$.

- Corollary 1.6.9: An $m \times n$ matrix is invertible if and only if its columns form a basis in \mathbb{F}^m .
- Linear span (of $\mathbf{v}_1, \dots, \mathbf{v}_n \in V$): The set of all linear combinations of $\mathbf{v}_1, \dots, \mathbf{v}_n$. Denoted by $\mathcal{L}\{\mathbf{v}_1, \dots, \mathbf{v}_n\}$, span $(\mathbf{v}_1, \dots, \mathbf{v}_n)$.

Systems of Linear Equations

2.1 Notes

9/29: • Row €

• Row elimination:

- Let

$$A = \begin{pmatrix} 1 & 2 & 3 & 1 \\ 3 & 1 & 2 & 7 \\ 2 & 1 & 2 & 1 \end{pmatrix}$$

- Then the **echelon form** matrix

$$A_e = \begin{pmatrix} 1 & 2 & 3 & 1 \\ 0 & 1 & 2 & -1 \\ 0 & 0 & 2 & -4 \end{pmatrix}$$

- Lastly, the **reduced echelon form** matrix

$$A_{re} = \begin{pmatrix} 1 & 0 & 0 & 7 \\ 0 & 1 & 0 & 3 \\ 0 & 0 & 1 & -2 \end{pmatrix}$$

• echelon form:

- All zero rows are below nonzero rows.

 For any nonzero row, its leading element is strictly to the left of the nonzero entry of the next row.

• Reduced echelon form:

– All pivots are 1.

- Used to solve systems of the form Ax = b.

• **Inconsistent** (system of equations): A system with no solution.

– If the last row is of the form $(0,\ldots,0,b)$ where $b\neq 0$, then there is no solution.

• Unique solution if A_e has a pivot in every column.

• There exists a solution for every b if there is a pivot in every row?

• Let $A: \mathbb{R}^n \to \mathbb{R}^m$ be a matrix. Then $\ker A = \{x \in \mathbb{R}^n : Ax = 0\}$ (subspace of \mathbb{R}^n) and range $A = \{Ax : x \in \mathbb{R}^n\}$ (subspace of \mathbb{R}^m).

• Also consider $\ker(A^T)$ and $\operatorname{range}(A^T)$, the basis of the kernel and range, and dimension.

- Finite-dimensional vector spaces:
 - A basis is a generating set (so every element of V can be written uniquely as a linear combination of the basis) the length of which is equal to the dimension of V.
 - All bases of finite-dimensional vector spaces have the same number of elements.
 - Let v_1, v_2, v_3 and w_1, w_2 be two generating sets of V.
 - Then

$$v_1 = \lambda_{11}w_1 + \lambda_{12}w_2$$

$$v_2 = \lambda_{21}w_1 + \lambda_{22}w_2$$

$$v_3 = \lambda_{31}w_1 + \lambda_{32}w_2$$

- Suppose the only solution to $\alpha_1v_1 + \alpha_2v_2 + \alpha_3v_3 = 0$ is $\alpha_1 = \alpha_2 = \alpha_3 = 0$.
- But this is not true, as we can find another one in terms of the λ s.
- If you have a list of linearly independent vectors, you can complete it into a basis.
 - If there exists a vector that can't be written as a linear combination of the list, add it to the list.
- If you find any particular solution to a system Ax = b, and you add to it any element of ker A, you will obtain another solution.
 - $Ax_1 = b$ and $Ax_h = 0$ implies that $A(x_1 + x_h) = b$.
 - $Ax_1 = b$ and $Ax_2 = b$ imply that $A(x_1 x_2) = 0$, i.e., that $x_1 x_2 \in \ker A$.
- If $A: \mathbb{R}^n \to \mathbb{R}^m$ and dim range A=m, then Ax=b is solveable for all $b \in \mathbb{R}^m$.
- Let rank $A = \dim \operatorname{range} A$.
- Rank theorem:
 - \blacksquare rank $A = \operatorname{rank} A^T$.
 - Let $A: \mathbb{R}^n \to \mathbb{R}^m$. We know that dim ker $A + \dim \operatorname{range} A = n$.
 - \blacksquare dim ker A^T + rank A^T = m.
 - This theorem survives linear algebra and enters functional analysis under the name **Fred-holm's alternative**.
- Fredholm's alternative: Ax = b has a solution for all $b \in \mathbb{R}^n$ iff dim ker $A^T = 0$.
 - dim ker $A^T = 0$ implies rank $A^T = m$ implies rank A = m implies dim range A = m, as desired.
- Pivot column (of A): A column of A where A_e has pivots.
- The **pivot columns** of A give a basis for range A.
- The pivot rows of A_e give a basis for range A^T .
- A basis for the kernel is enough to solve Ax = 0.
- If you take these three things as givens, you can prove the rank theorem.

2.2 Chapter 2: Systems of Linear Equations

From Treil (2017).

10/24:

- A system is inconsistent iff the echelon form of the augmented matrix has a row of the form $(0 \cdots 0 \ b)$.
 - A solution to $A\mathbf{x} = \mathbf{b}$ is unique iff there are no free variables, i.e., iff there is a pivot in every column.
 - $A\mathbf{x} = \mathbf{b}$ is consistent iff the echelon form of the coefficient matrix has a pivot in every row.

- $A\mathbf{x} = \mathbf{b}$ has a unique solution for any \mathbf{b} iff the echelon form of the coefficient matrix A has a pivot in every row and column.
- Proposition 2.3.1: Let $\mathbf{v}_1, \dots, \mathbf{v}_m \in \mathbb{F}^n$, and let $A = \begin{bmatrix} \mathbf{v}_1 & \cdots & \mathbf{v}_m \end{bmatrix}$ be an $n \times m$ matrix with columns $\mathbf{v}_1, \dots, \mathbf{v}_m$. Then
 - 1. The system $\mathbf{v}_1, \dots, \mathbf{v}_m$ is linearly independent iff the echelon form of A has a pivot in every column.
 - 2. The system $\mathbf{v}_1, \dots, \mathbf{v}_m$ is complete iff the echelon form of A has a pivot in every row.
 - 3. The system $\mathbf{v}_1, \dots, \mathbf{v}_m$ is a basis of \mathbb{F}^n iff the echelon form of A has a pivot in every column and in every row.
- Proposition 2.3.6: A matrix A is invertible if and only if its echelon form has a pivot in every column and every row.
- Corollary 2.3.7: An invertible matrix must be square $(n \times n)$.
- Proposition 2.3.8: If a square $(n \times n)$ matrix is left invertible or if it is right invertible, then it is invertible. In other words, to check the invertibility of a square matrix A, it is sufficient to check only one of the conditions $AA^{-1} = I$, $A^{-1}A = I$.
- Any invertible matrix is row-equivalent to (can be row-reduced to) to the identity matrix.
- Homogeneous (system of linear equations): A system of the form $A\mathbf{x} = \mathbf{0}$.
- Theorem 2.6.1: Let a vector \mathbf{x}_1 satisfy the equation $A\mathbf{x} = \mathbf{b}$. and let H be the set of all solutions of the associated homogeneous system $A\mathbf{x} = \mathbf{0}$. Then the set

$$\{\mathbf{x}_1 + \mathbf{x}_h : \mathbf{x}_h \in H\}$$

is the set of all solutions to the equation $A\mathbf{x} = \mathbf{b}$.

- The pivot columns are a basis of range A. The pivot rows are a basis of range A^T . The solutions to the equation $A\mathbf{x} = \mathbf{0}$ are a basis of ker A.
- Theorem 2.7.3: Let A be an $m \times n$ matrix. Then the equation $A\mathbf{x} = \mathbf{b}$ has a solution for every $\mathbf{b} \in \mathbb{R}^m$ iff the dual equation $A^T\mathbf{x} = \mathbf{0}$ has a unique (only the trivial) solution.
 - Note that this is a corollary to the rank theorem.
- Change of coordinates formula:
 - Let $T: V \to W$ be a linear transformation, and let $\mathcal{V} = \{\mathbf{v}_1, \dots, \mathbf{v}_n\}$ and $\mathcal{W} = \{\mathbf{w}_1, \dots, \mathbf{w}_n\}$ be bases of V and W, respectively.
 - The $m \times n$ matrix of T with respect to these bases is $[T]_{WV}$, and relates the coordinates of $[T\mathbf{v}]_{W}$ and $[\mathbf{v}]_{V}$ via

$$[T\mathbf{v}]_{\mathcal{W}} = [T]_{\mathcal{W}\mathcal{V}}[\mathbf{v}]_{\mathcal{V}}$$

- Change of coordinates matrix: If \mathcal{A}, \mathcal{B} are two bases of V, then we can convert the coordinates of a vector in \mathcal{B} to its in \mathcal{A} with the identity matrix (with respect to the appropriate bases). In particular,

$$[\mathbf{v}]_{\mathcal{B}} = [I]_{\mathcal{B}\mathcal{A}}[\mathbf{v}]_{\mathcal{A}}$$

- Note that the k^{th} column of $[I]_{\mathcal{BA}}$ is the coordinate representation in \mathcal{B} of \mathbf{a}_k , i.e., $[\mathbf{a}_k]_{\mathcal{B}}$.
- The change of coordinates matrix from a basis \mathcal{B} to the standard basis \mathcal{S} is easy to compute; by the above, it's just

$$[I]_{\mathcal{SB}} = \begin{bmatrix} \mathbf{b}_1 & \cdots & \mathbf{b}_n \end{bmatrix}$$

■ It follows that $[I]_{\mathcal{BS}} = ([I]_{\mathcal{SB}})^{-1}$.

- \blacksquare This allows us to compute $[I]_{\mathcal{BA}}$ as $[I]_{\mathcal{BS}}[I]_{\mathcal{SA}}$
- If $T: V \to W$, $\mathcal{A}, \tilde{\mathcal{A}}$ are bases of V, and $\mathcal{B}, \tilde{\mathcal{B}}$ are bases of W, and we have $[T]_{\mathcal{B}\mathcal{A}}$, then

$$[T]_{\tilde{\mathcal{B}}\tilde{\mathcal{A}}} = [I]_{\tilde{\mathcal{B}}\mathcal{B}}[T]_{\mathcal{B}\mathcal{A}}[I]_{\mathcal{A}\tilde{\mathcal{A}}}$$

• Change of basis ends up at similarity; two operators are similar if we can change the basis of one into another.

Determinants

3.1 Notes

9/29: • The determinant, geometrically, is the volume of the object (in \mathbb{R}^3) you get when you take linear combinations of the vectors.

• In 2D:

10/1:

- Let v_1, v_2 be two vectors. Put tail to tail and forming a parallelogram, the determinant of the matrix (v_1, v_2) is the area of said parallelogram.
- Linearity 1: $D(av_1, v_2, ..., v_n) = aD(v_1, ..., v_n)$ is the same as saying that if you stretch one vector by a, you scale up the area by that much, too.
- Linearity 2: $D(v_1, \ldots, v_{k+} + v_{k-}, \ldots, v_n) = D(-) + D(+)$.
- Antisymmetry: $D(v_1, \ldots, v_k, \ldots, v_j, \ldots, v_n) = -D(v_1, \ldots, v_j, \ldots, v_k, \ldots, v_n)$. Interchanging columns flips the sign of the determinant.
- Basis: $D(e_1, ..., e_n) = 1$.
- Determinant: Denoted by $D(v_1, \ldots, v_n)$, where (v_1, \ldots, v_n) is an $n \times n$ matrix.
- Consider an $n \times n$ matrix A consisting of n columns containing vectors $\mathbf{v}_1, \dots, \mathbf{v}_n \in \mathbb{R}^n$.
 - -D(A) is the volume of the solid $V = \sum_{i=1}^{n} \alpha_i v_i$.
 - $D(\mathbf{e}_1, \dots, \mathbf{e}_n) = 1.$

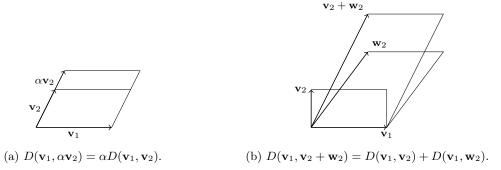


Figure 3.1: Visualizing properties of determinants.

- Basic properties of the determinant.
 - If A has a zero column, then $\det A = 0$: Scalar property.

- If A has two equal columns, then $\det A = 0$: Multiply one by minus and add.
- If A has a column which is a multiple of another, then $\det A = 0$: Pull out the multiple and then you have the previous one.
- If columns are linearly dependent, then $\det A = 0$: Decompose it into sums, split, add back up with previous properties.
- The determinant is preserved under column reduction.
- $-\det A^T = \det A$: Put everything in rref.
- If A is not invertible, then $\det A = 0$ (not invertible implies linearly dependent columns, implies $\det A = 0$).
- $-\det(AB) = \det A \det B.$
- Determinant of...
 - A diagonal matrix: The product of the diagonal entries (pull out the terms, and then note that the remaining identity matrix has determinant 1).
 - An upper triangular matrix: The product of the diagonal entries (column reduction to make it into a diagonal matrix, and then the property above).

3.2 Chapter 3: Determinants

From Treil (2017).

10/24:

- Let $A_{j,k}$ denote the $(n-1) \times (n-1)$ matrix obtained from A by crossing out row j and column k and pushing it together.
- Cofactors (of A): The numbers $C_{j,k}$, one per entry, defined by

$$C_{j,k} = (-1)^{j+k} \det A_{j,k}$$

• Cofactor matrix (of A): The matrix

$$C = \{C_{j,k}\}_{j,k=1}^{n}$$

• Theorem 3.5.2: Let A be an invertible matrix and let C be its cofactor matrix. Then

$$A^{-1} = \frac{1}{\det A} C^T$$

• Cramer's rule: If A is invertible and $A\mathbf{x} = \mathbf{b}$, then

$$x_k = \frac{\det B_k}{\det A}$$

where B_k is obtained from A by replacing column k of A by the vector **b**.

- Minor (of order k of A): The determinant of a $k \times k$ submatrix of A.
- Theorem 3.6.1: The rank of a nonzero matrix A is equal to the largest integer k such that there exists a nonzero minor of order k.

Introduction to Spectral Theory

4.1 Notes

- **Difference equation**: Like a differential equation, but instead of writing a differentials, you write differences.
 - Suppose we want to solve $x_{n+1} = Ax_n$ with x_0 given.
 - You will find that $x_n = A^n x_0$.
 - This gets hard to compute, so we want to find a way to simplify the computation.
 - Thus, we want to diagonalize the matrix, and this concept is inherently linked to eigenvalues and eigenvectors.
 - If you can decompose the x_0 into a linear combination of eigenvectors, then you can simplify the computation a lot:

$$x_n = \sum \alpha_i A^n v_i = \sum \alpha_i \lambda_i^n v_i$$

- An $n \times n$ matrix will have n eigenvalues. You want n linearly independent eigenvectors, creating an eigenbasis.
- To find eigenvalues and eigenvectors, we need to solve $Ax = \lambda x$, i.e., $(A \lambda I)x = 0$. Thus, $\ker(A \lambda I) \neq \{0\}$, so $\det(A \lambda I) = 0$.
- The eigenvalues of A are independent of the choice of basis of the domain of A or the range.
- We need to know everything in Treil (2017).
 - We don't need to know the applications sections, but you should be interested.
 - Spectral theory: Decomposing a linear operator.
 - Let $A:V\to V$ be a linear operator. $\lambda\in\mathbb{C}$ is an eigenvalue if there exists $x\in V$ nonzero such that $Ax=\lambda x$.
 - Let A be an $n \times n$ matrix over \mathbb{C} or \mathbb{R} .
 - The eigenvalues are the roots of the polynomial $det(A \lambda I) = 0$ in λ .
 - Things we want to do:
 - Given A, find the eigenvalues and eigenvectors (solve $(A \lambda I)x = 0$).

- In order to simplify A, make it a diagonal matrix:

$$A = S \begin{pmatrix} \lambda_1 & & 0 \\ & \ddots & \\ 0 & & \lambda_n \end{pmatrix} S^{-1}$$

- Eigenvalues are independent of the choice of basis.
 - From the book, we have that

$$[A]_{\mathcal{A}\mathcal{A}} = [S]_{\mathcal{A}\mathcal{B}}[B]_{\mathcal{B}\mathcal{B}}[S]_{\mathcal{A}\mathcal{B}}^{-1}$$

- It follows that

$$A - \lambda I = [S]_{\mathcal{AB}}(B - \lambda I)[S]_{\mathcal{AB}}^{-1}$$

SO

$$\det(A - \lambda I) = \det([S]_{\mathcal{AB}}(B - \lambda I)[S]_{\mathcal{AB}}^{-1}) = \det([S]_{\mathcal{AB}}[S]_{\mathcal{AB}}^{-1}(B - \lambda I)) = \det(B - \lambda I)$$

- If $p(z) = (z \lambda)^k q(z)$, then k is the algebraic multiplicity of λ . The geometric multiplicity of λ is dim $\ker(A \lambda I)$.
 - These terms are not always the same, but they are related.
- Diagonalization:
 - Given A that corresponds to $T:V\to V$, can we find a basis of V in which the operator is a diagonal matrix?
 - $-A = SDS^{-1}$ iff there exists a basis of V consisting of the eigenvectors of A.
 - Proves $A^{N} = SD^{N}S^{-1}$ via $A^{2} = SDS^{-1}SDS^{-1} = SDIDS^{-1} = SD^{2}S^{-1}$.
- Let A be an $n \times n$ matrix over \mathbb{F} . If $\lambda_1, \ldots, \lambda_r$ are distinct eigenvalues, then their eigenvectors are linearly independent.
 - Prove with induction contradiction argument. Assume true for \mathbf{v}_{r-1} . Then

$$0 = (A - \lambda_r I)[\mathbf{v}_1 + \dots + \mathbf{v}_r] = (\lambda_1 - \lambda_r)\mathbf{v}_1 + \dots + (\lambda_{r-1} - \lambda_r)\mathbf{v}_{r-1}$$

- Implies $\lambda_r = \lambda_i$ for all $i \in [r-1]$, a contradiction.
- If A has n distinct eigenvalues, then A is diagonalizable.
- If $A: V \to V$ has n complex eigenvalues, then A is diagonalizable iff the algebraic multiplicity equals the geometric multiplicity for each eigenvalue.
- Goes through a sample diagonalization with $\begin{pmatrix} 1 & 2 \\ 8 & 1 \end{pmatrix}$.
 - We have

$$A - \lambda I = \begin{pmatrix} 1 - \lambda & 2 \\ 8 & 1 - \lambda \end{pmatrix}$$

so

$$0 = \det(A - \lambda I) = (1 - \lambda)^2 - 16$$

- It follows that $\lambda = 5, -3$.
- This yields

$$\begin{pmatrix} 1 & 2 \\ 8 & 1 \end{pmatrix} = \begin{pmatrix} 1 & 1 \\ 2 & -2 \end{pmatrix} \begin{pmatrix} 5 & 0 \\ 0 & -3 \end{pmatrix} \begin{pmatrix} 1 & 1 \\ 2 & -2 \end{pmatrix}^{-1}$$

by inspection.

- As another example, consider $\begin{pmatrix} 1 & 2 \\ -2 & 1 \end{pmatrix}$.
 - Here, we have $\lambda = 1 \pm 2i$.

4.2 Chapter 4: Introduction to Spectral Theory

From Treil (2017).

10/24:

- **Spectrum** (of A): The set of all eigenvalues of A. Denoted by $\sigma(A)$.
- Proposition 4.1.1: The geometric multiplicity of an eigenvalue cannot exceed its algebraic multiplicity.
- Theorem 4.2.1: A matrix A (with values in \mathbb{F}) admits a representation $A = SDS^{-1}$ where D is a diagonal matrix and S is invertible if and only if there exists a basis of \mathbb{F}^n of eigenvectors of A. Moreover, in this case diagonal entries of D are the eigenvalues of A and columns of S are the corresponding eigenvectors.
- Any operator on a complex vector space has n eigenvalues (counting multiplicities).
 - Think n necessary roots of the characteristic polynomial, or the necessary upper triangular representation.
- Theorem 4.2.8: Let an operator $A: V \to V$ have exactly $n = \dim V$ eigenvalues (counting multiplicities). Then A is diagonalizable if and only if for each eigenvalue λ , the dimension of the eigenspace $\ker(A \lambda I)$ (i.e., the geometric multiplicity of λ) coincides with the algebraic multiplicity of λ .
- Theorem 4.2.9: A real $n \times n$ matrix A admits a real factorization (i.e., a real representation $A = SDS^{-1}$ where S and D are real matrices, D is diagonal, and S is invertible) if and only if it admits a complex factorization and all eigenvalues of A are real.
- Example of a nondiagonalizable matrix:

$$A = \begin{pmatrix} 1 & 1 \\ 0 & 1 \end{pmatrix}$$

- $-p(\lambda)=(1-\lambda)^2$, so $\lambda=1$ with algebraic multiplicity 2.
- However, dim $\ker(A-I) = 1$ since A-I has only one pivot, hence 2-1=1 free variable.
- Thus, apply Theorem 4.2.8.

Inner Product Spaces

5.1 Notes

10/6:

• We define

$$\ell^{2}(\mathbb{R}) = \left\{ \{a_{n}\}_{n \geq 1} \subset \mathbb{R} : \sum_{1}^{\infty} |a_{n}|^{2} < \infty \right\}$$

- Inner product: A map $V \times V \to \mathbb{F}$ that takes $(\mathbf{x}, \mathbf{y}) \mapsto \mathbf{x} \cdot \mathbf{y}$. Denoted by $\cdot, (\cdot, \cdot), \langle \cdot, \cdot \rangle$.
- Properties of the inner product:

$$-(\mathbf{x}, \mathbf{y}) = \overline{(\mathbf{y}, \mathbf{x})}$$
 (symmetry).

$$- (\alpha \mathbf{x} + \beta \mathbf{y}, \mathbf{z}) = \alpha(\mathbf{x}, \mathbf{z}) + \beta(\mathbf{y}, \mathbf{z}) \text{ (linearity)}.$$

$$- (\mathbf{x}, \mathbf{x}) \ge 0.$$

$$- (\mathbf{x}, \mathbf{x}) = 0 \text{ iff } \mathbf{x} = 0.$$

• If $\mathbf{x}, \mathbf{y} \in \mathbb{R}^n$, then

$$(\mathbf{x}, \mathbf{y}) = \sum_{i=1}^{n} x_i y_i$$

• If $\mathbf{x}, \mathbf{y} \in \mathbb{C}^n$, then

$$(\mathbf{x}, \mathbf{y}) = \sum_{i=1}^{n} x_i \bar{y}_i$$

• If $f, g \in \mathbb{P}_n(t)$, then

$$(f,g) = \int_{-1}^{1} f\bar{g} \,\mathrm{d}t$$

- The conjugate of a polynomial is the polynomial with the conjugate of the coefficients of the original polynomial. Symbolically, if $f = \sum_{i=0}^{n} \alpha_i t^i$ is a polynomial, then $\bar{f} = \sum_{i=0}^{n} \bar{\alpha}_i t^i$.
- It is a fact that

$$\left| \sum_{n=0}^{\infty} a_n \bar{b}_n \right| \le \| (a_n)_{n \ge 1} \| \| (b_n)_{n \ge 1} \|$$

- Suppose we want to define the inner product between two matrices.
 - A common one is

$$(A, B) = \operatorname{tr}(B^*A)$$

where $B^* = \overline{B}^T = \overline{B^T}$ is the conjugate transpose.

• We define the norm as a function $V \to [0, \infty)$ given by

$$\|\mathbf{x}\| = \sqrt{(\mathbf{x}, \mathbf{x})}$$

- Properties of the norm.
 - $\|\alpha \mathbf{x}\| = |\alpha| \|\mathbf{x}\|.$
 - $\|\mathbf{x} + \mathbf{y}\| \le \|\mathbf{x}\| + \|\mathbf{y}\|.$
 - $\|\mathbf{x}\| = 0 \text{ iff } \mathbf{x} = 0.$
- In \mathbb{R}^n ,



Figure 5.1: The unit ball of norms corresponding to $p = 1, 2, \infty$.

- The standard norm is

$$\|\mathbf{x}\| = \sqrt{\sum |x_i|^2}$$

- We can also define

$$\|\mathbf{x}\|_p = \sqrt[p]{\sum |x_i|^p}$$

- We can even define

$$\|\mathbf{x}\|_{\infty} = \max|x_i|$$

- And we can prove that all of these are valid norms.
- Only the norm corresponding to ℓ^2 is given by an inner product, but all the other quantities are still norms as defined by the properties (see Treil (2017)).
- Figure 5.1 shows the unit ball of each norm, i.e., the set of all points which have norm 1.
- The parallelogram rule:

$$\|\mathbf{x} + \mathbf{y}\|^2 + \|\mathbf{x} - \mathbf{y}\|^2 = 2(\|\mathbf{x}\|^2 + \|\mathbf{y}\|^2)$$

- Orthogonality: Given \mathbf{v}, \mathbf{w} , if $\mathbf{v} \perp \mathbf{w}$, then $(\mathbf{v}, \mathbf{w}) = 0$.
- In particular, if $\mathbf{v} \perp \mathbf{w}$, then

$$\|\mathbf{v} + \mathbf{w}\|^2 = \|\mathbf{v}\|^2 + \|\mathbf{w}\|^2$$

- Let E be a subspace of V. If $\mathbf{v} \perp E$, then $\mathbf{v} \perp \mathbf{e}$ for all $\mathbf{e} \in E$, i.e., $\mathbf{v} \perp \mathbf{a}$ set of vectors spanning E.
- Any set of orthogonal vectors is linearly independent. Thus, if V is n dimensional, then $\mathbf{v}_1, \dots, \mathbf{v}_n$ orthogonal is a basis.
- Let E be a subspace of V. Take $\mathbf{v} \in V$. We want to define the projection $P_E \mathbf{v}$ of \mathbf{v} onto E.
 - We have that $P_E \mathbf{v} \in E$ and $v P_E \mathbf{v} \perp E$.
 - Additionally, we have that

$$\|\mathbf{v} - P_E \mathbf{v}\| \le \|\mathbf{v} - \mathbf{e}\|$$

for all $\mathbf{e} \in E$.

- Lastly, we have that $P_E \mathbf{v}$ is unique.
- If we receive a basis of a vector space, how do we create out of that a basis that is orthogonal? The process of doing this is called **Gram-Schmidt orthogonalization**.
 - We keep \mathbf{v}_1 , subtract $P_{\mathbf{v}_1}\mathbf{v}_2$ from \mathbf{v}_2 , subtract $P_{\{\mathbf{v}_1,\mathbf{v}_2\}}\mathbf{v}_3$ from \mathbf{v}_3 , and on and on.
- If we are given a set of orthogonal vectors, we can normalize them by dividing each by its norm. This creates an orthonormal list. The standard basis is orthonormal.
- Let

$$E^{\perp} = \{ v \in V : v \perp E \}$$

- It follows that $V = E \oplus E^{\perp}$.
- How close can we come to solving $A\mathbf{x} = \mathbf{b}$ if we cannot solve it exactly (i.e., if the columns are not linearly independent)?
 - Let A be an $m \times n$ matrix, and let $\mathbf{b} \in \mathbb{R}^m$.
 - Then the best solution is given by minimizing $||A\mathbf{x} \mathbf{b}||$. We minimize this with projections. A special case of this is least squares regression! More details in Treil (2017).
- 10/8: Soug is gonna send us a hefty amount of reading for the weekend.
 - Least square approximation:
 - If we want to minimize $||A\mathbf{x} \mathbf{b}||$, the best we can do is project **b** onto the range of A.
 - Let $\mathbf{v}_1, \dots, \mathbf{v}_k$ be an orthogonal basis of range A.
 - Then

$$\operatorname{Proj}_{\operatorname{range} A} \mathbf{b} = \sum_{k=1}^{k} \frac{(\mathbf{b}, \mathbf{v}_{k})}{\|v_{k}\|^{2}} \mathbf{v}_{k}$$

- Matrix equation form:

$$Projection_{range A} = A(A^*A)^{-1}A^*$$

if A^*A is invertible, where $A^* = \bar{A}^T$.

- Soug never uses this though.
- The minimum is found when $\mathbf{b} A\mathbf{x} \perp \text{range } A$. Implies that $\mathbf{b} A\mathbf{x} \perp \mathbf{a}_k$ for all k. Implies $(\mathbf{b} A\mathbf{x}, \mathbf{a}_k) = \bar{\mathbf{a}}_k^T(\mathbf{b} A\mathbf{x}) = 0$.
- Note that we're letting $\bar{\mathbf{a}}_k^T$ be the row vector

$$\bar{\mathbf{a}}_k^T = \begin{pmatrix} \bar{a}_{1,k} & \cdots & \bar{a}_{n,k} \end{pmatrix}$$

- We also have $\bar{A}^T(\mathbf{b} A\mathbf{x}) = 0$, from which it follows that $A^*A\mathbf{x} = A^*\mathbf{b}$, so $\mathbf{x} = (A^*A)^{-1}A^*\mathbf{b}$. Thus, $\text{Proj}|_{\text{range }A} = Ax$, so $\text{Proj}|_{\text{range }A} = A(A^*A)^{-1}A^*\mathbf{b}$.
- Adjoint of a linear map $T: V \to W$ is the A^* discussed above.
 - First, we'll do this for matrices. And then we'll do it for any finite-dimensional vector space.
 - Let A be an $m \times n$ matrix. We claim then that

$$(A\mathbf{x}, \mathbf{y}) = (\mathbf{x}, A^*\mathbf{y})$$

for all $\mathbf{x} \in \mathbb{C}^n$, $\mathbf{y} \in \mathbb{C}^m$. Proof:

$$(A\mathbf{x}, \mathbf{y}) = \bar{\mathbf{y}}^T A \mathbf{x}$$
$$= \mathbf{y}^* A \mathbf{x}$$
$$= (A^* \mathbf{y})^* \mathbf{x}$$
$$= (\mathbf{x}, A^* \mathbf{y})$$

- Properties of the adjoint:

$$(AB)^{T} = B^{T}A^{T}$$
$$(AB)^{*} = B^{*}A^{*}$$
$$(A^{*})^{*} = A$$

- $-A^*$ is the unique matrix B such that $(A\mathbf{x}, \mathbf{y}) = (\mathbf{x}, B\mathbf{y})$.
- Let $\mathbf{v}_1, \ldots, \mathbf{v}_n$ be a basis of V, and let $\mathbf{w}_1, \ldots, \mathbf{w}_m$ be a basis of W.
- Definition of A^* : If $(A\mathbf{x}, \mathbf{y}) = (y, A^*\mathbf{x})$ for all $\mathbf{x} \in V$ and $\mathbf{y} \in W$.
- But it's not enough to define something; we have to check that it exists.
- If $[A]_{\mathcal{AB}}$, then $[A^*]_{\mathcal{AB}}$.
- More properties (give criteria for solving systems of equations):

$$\ker A^* = (\operatorname{range} A)^{\perp}$$

 $\ker A = (\operatorname{range} A^*)^{\perp}$
 $\operatorname{range} A = (\ker A^*)^{\perp}$
 $\operatorname{range} A^* = (\ker A)^{\perp}$

- Soug proves these.
- Isometries and unitary operators.
 - $-U: X \to Y$ is an isometry if $\|\mathbf{x}\| = \|U\mathbf{x}\|$ for all $\mathbf{x} \in X$. It is an isometry because it preserves the distance between points.
 - It immediately follows that $\|\mathbf{x}_1 \mathbf{x}_2\| = \|U\mathbf{x}_1 U\mathbf{x}_2\| = \|U(\mathbf{x}_1 \mathbf{x}_2)\|$.
 - This definition is equivalent to an inner product one: $(\mathbf{x}, \mathbf{y}) = (U\mathbf{x}, U\mathbf{y})$. This follows from the definition of the norm.
 - We have

$$(\mathbf{a}, \mathbf{b}) = \frac{1}{4} \sum_{\alpha = \pm 1, \pm i} \alpha \|\mathbf{a} + \alpha \mathbf{b}\|^2$$

■ $(a+b)^2 - (a-b)^2 = 4ab$ for any $a, b \in \mathbb{R}$, so $ab = \frac{1}{4}[(a+b)^2 - (a-b)^2]$. Thus, in a real inner product space,

$$(\mathbf{a}, \mathbf{b}) = \frac{1}{4} (\|\mathbf{a} + \mathbf{b}\|^2 - \|\mathbf{a} - \mathbf{b}\|^2)$$

- It follows that isometries preserve inner products.
- U is an isometry if and only if $U^*U = I$. Proof:

$$(\mathbf{x}, \mathbf{x}) = (U\mathbf{x}, U\mathbf{y}) = (U^*U\mathbf{x}, \mathbf{x})$$
$$(\mathbf{x}, \mathbf{y}) = (U\mathbf{x}, U\mathbf{y}) = (U^*U\mathbf{x}, \mathbf{y})$$
$$(\mathbf{x}, \mathbf{y}) = (U\mathbf{x}, U\mathbf{y}) = (\mathbf{x}, \mathbf{y})$$

for all \mathbf{y} .

- An isometry is unitary if it is invertible.
 - Thus, $U: X \to Y$ an isometry is unitary iff dim $X = \dim Y$.
- Note that it follows that $U^* = U^{-1}$ for U an isometry.
- U unitary implies $|\det U| = 1$, so λ an eigenvalue of U implies that $|\lambda| = 1$.
- A is diagonalizable iff it has an orthogonal basis of eigenvectors.

5.2 Chapter 5: Inner Product Spaces

From Treil (2017).

10/24:

• Standard inner product (on \mathbb{C}^n): The inner product (\mathbf{z}, \mathbf{w}) defined by

$$(\mathbf{z}, \mathbf{w}) = \mathbf{w}^* \mathbf{z}$$

• Corollary 5.1.5: Let \mathbf{x}, \mathbf{y} be vectors in an inner product space V. The equality $\mathbf{x} = \mathbf{y}$ holds if and only if

$$(\mathbf{x}, \mathbf{z}) = (\mathbf{y}, \mathbf{z})$$

for all $\mathbf{z} \in V$.

• Corollary 5.1.6: Suppose two operator $A, B: X \to Y$ satisfy

$$(A\mathbf{x}, \mathbf{y}) = (B\mathbf{x}, \mathbf{y})$$

for all $\mathbf{x} \in x$ and $\mathbf{y} \in Y$. Then A = B.

- Normed space: A vector space V equipped with a norm that satisfies properties of homogeneity, the triangle inequality, non-negativity, and non-degeneracy.
- Any inner product space is naturally a normed space.
- If $1 \leq p < \infty$, we can define a corresponding norm on \mathbb{R}^n or \mathbb{C}^n by

$$\|\mathbf{x}\|_p = \left(\sum_{k=1}^n |x_k|^p\right)^{1/p}$$

• We can also define the norm for $p = \infty$ by

$$\|\mathbf{x}\|_{\infty} = \max\{|x_k| : k = 1, \dots, n\}$$

- Note that the norm of this form for p=2 is the usual norm.
- These norms are heavily associated with Figure 5.1.
- Minkowski inequality: One of the triangle inequalities for norms with $p \neq 2$.
- Theorem 5.1.11: A norm in a normed space is obtained from some inner product if and only if it satisfies the Parallelogram Identity

$$\|\mathbf{u} + \mathbf{v}\|^2 + \|\mathbf{u} - \mathbf{v}\|^2 = 2(\|\mathbf{u}\|^2 + \|\mathbf{v}\|^2)$$

for all $\mathbf{u}, \mathbf{v} \in V$.

- It follows that norms with $p \neq 2$ do not have associated inner products, since such norms fail to satisfy the parallelogram identity.
- Lemma 5.2.5 (Generalized Pythagorean Identity): Let $\mathbf{v}_1, \dots, \mathbf{v}_n$ be an orthogonal system. Then

$$\left\| \sum_{k=1}^{n} \alpha_k \mathbf{v}_k \right\|^2 = \sum_{k=1}^{n} |\alpha_k|^2 \|\mathbf{v}_k\|^2$$

• Proposition 5.3.3: Let $\mathbf{v}_1, \dots, \mathbf{v}_r$ be an orthogonal basis in E. Then the orthogonal projection $P_E \mathbf{v}$ of a vector \mathbf{v} is given by the formula

$$P_E \mathbf{v} = \sum_{k=1}^r \frac{(\mathbf{v}, \mathbf{v}_k)}{\|\mathbf{v}_k\|^2} \mathbf{v}_k$$

- It follows that

$$P_{E}\mathbf{v} = \sum_{k=1}^{r} \frac{\mathbf{v}_{k}^{*}\mathbf{v}}{\|\mathbf{v}_{k}\|^{2}} \mathbf{v}_{k}$$

$$= \sum_{k=1}^{r} \frac{1}{\|\mathbf{v}_{k}\|^{2}} \mathbf{v}_{k} \mathbf{v}_{k}^{*} \mathbf{v}$$

$$= \left(\sum_{k=1}^{r} \frac{1}{\|\mathbf{v}_{k}\|^{2}} \mathbf{v}_{k} \mathbf{v}_{k}^{*}\right) \mathbf{v}$$

- Thus, we have that

$$P_E = \sum_{k=1}^r \frac{1}{\|\mathbf{v}_k\|^2} \mathbf{v}_k \mathbf{v}_k^*$$

- Gram-Schmidt orthogonalization: Let $\mathbf{x}_1, \dots, \mathbf{x}_n$ be a linearly independent system of vectors to orthogonalize. Then $\mathbf{v}_1 = \mathbf{x}_1$, $\mathbf{v}_2 = \mathbf{x}_2 P_{\text{span}\{\mathbf{v}_1\}}\mathbf{x}_2$, $\mathbf{v}_3 = \mathbf{x}_3 P_{\text{span}\{\mathbf{v}_1,\mathbf{v}_2\}}\mathbf{x}_3$, and on and on.
- To find the least squares solution to $A\mathbf{x} = \mathbf{b}$, solve $A\mathbf{x} = P_{\text{range }A}\mathbf{b}$.
 - We can do this by finding an orthogonal basis of range A and then applying the projection formula.
 - Alternatively, we can use the following formula to speed things up if A^*A is invertible:

$$P_{\text{range }A}\mathbf{b} = A(A^*A)^{-1}A^*\mathbf{b}$$

• Theorem 5.4.1: For an $m \times n$ matrix A,

$$\ker A = \ker(A^*A)$$

- Thus, A^*A is invertible iff A is invertible iff A is full rank. This gives us a condition on when we can use the projection formula.
- Theorem 5.6.1: An operator $U: X \to Y$ is an isometry if and only if it preserves the inner product, i.e., if and only if

$$(\mathbf{x}, \mathbf{y}) = (U\mathbf{x}, U\mathbf{y})$$

for all $\mathbf{x}, \mathbf{y} \in X$.

- Lemma 5.6.2: An operator $U: X \to Y$ is an isometry if and only if $U^*U = I$.
- Unitary (operator): An invertible isometry.
- Proposition 5.6.3: An isometry $U: X \to Y$ is a unitary operator iff $\dim X = \dim Y$.
- Orthogonal (matrix): A unitary matrix with real entries.
- Unitary operator properties:
 - 1. $U^{-1} = U^*$.
 - 2. U unitary implies $U^* = U^{-1}$ unitary.
 - 3. If $\mathbf{v}_1, \dots, \mathbf{v}_n$ is orthonormal, $U\mathbf{v}_1, \dots, U\mathbf{v}_n$ is orthonormal.
 - 4. U, V unitary implies UV unitary.
- \bullet A matrix U is an isometry iff its columns form an orthonormal system.
- \bullet Proposition 5.6.4: Let U be a unitary matrix. Then
 - 1. $|\det U| = 1$. In particular, if U is orthogonal, then $\det U = \pm 1$.
 - 2. $|\lambda| = 1$ for every eigenvalue λ of U.
- Proposition 5.6.5: A matrix A is unitarily equivalent to a diagonal one iff it has an orthogonal (orthonormal) basis of eigenvectors.

Structure of Operators on Inner Product Spaces

6.1 Notes

- 10/11: Spectral decomposition of self-adjoint linear maps.
 - Can we write a map in term of the eigenvalues only?
 - Let $A: X \to X$ be linear and self-adjoint. Where dim $X < \infty$.
 - Let A have eigenvalues $\lambda_1, \ldots, \lambda_n$ and eigenvectors $\mathbf{v}_1, \ldots, \mathbf{v}_n$. The there is an orthonormal basis of X consisting of eigenvectors of A. An operator is self-adjoint if $A = A^*$.
 - If A is self-adjoint, then A can be written as diagonal with the eigenvalues on the diagonal with respect to some orthonormal basis of eigenvectors.
 - Let $\mathbb{F} = \mathbb{C}$.
 - If there exists an orthonormal basis u_1, \ldots, u_n of X such that A is triangular, then $A = UTU^*$ where U is unitary and T is upper triangular.
 - Proved with induction on dim X.
 - $-\dim X = 1$ is clear.
 - Assume for dim X = n 1, WTS for dim X = n.
 - The subspace has a basis $\mathbf{v}_1, \dots, \mathbf{v}_{n-1}$ such that A has a diagonal form.
 - Let $u \in X$ be linearly independent of $\mathbf{v}_1, \dots, \mathbf{v}_{n-1}$.
 - Let λ be the remaining eigenvalue and u the corresponding eigenvector. Let E = span(u). Then make the matrix λ in the upper left corner, and block diagonal with " A_{n-1} " in the bottom right corner, zeroes everywhere else.
 - Self-adjoint (matrix A): A linear map $A: X \to X$ where dim $X < \infty$ such that $A = A^*$.
 - Similarly, (Ax, y) = (x, Ay).
 - A self-adjoint implies all eigenvalues are real, eigenvectors corresponding to different eigenvalues are orthogonal.
 - Soug proves this.
 - Strictly positive (operator A): A self-adjoint operator $A: X \to X$ such that (Ax, x) > 0 for all $x \neq 0$. Also known as positive definite.
 - Implies that all eigenvalues are strictly positive.

- Nonnegative (operator A): A self-adjoint operator $A: X \to X$ such that $(Ax, x) \ge 0$ for all $x \ne 0$. Also known as definite.
 - All eigenvalues are nonnegative.
- Suppose $A \ge 0$ is self-adjoint. Then there exists a unique self-adjoint $B \ge 0$ such that $B^2 = A$.
 - A self-adjoint is diagonal (wrt. some basis).
 - A positive means that all eigenvalues (diagonal entries) are positive.
 - Thus, take

$$B = \begin{pmatrix} \sqrt{\lambda_1} & 0 & 0 \\ 0 & \ddots & 0 \\ 0 & 0 & \sqrt{\lambda_n} \end{pmatrix}$$

- Suppose $B^2 = A$, $C^2 = A$. Then we have an orthonormal basis corresponding to B and an orthonormal basis corresponding to C. It follows that $B^2 = C^2 = A$. Write B^2x and C^2x in terms of their bases; will necessitate that the bases are the same.
- 10/13: If we get yes/no questions, we don't have to justify.
 - Cauchy-Schwarz inequality:

$$|(\mathbf{x}, \mathbf{y})| \le \|\mathbf{x}\| \|\mathbf{y}\|$$

- Real spaces, V vs. (\cdot, \cdot) inner product.
- Proof:

$$0 \le \|\mathbf{x} + t\mathbf{y}\|^2$$
$$= t^2 \|\mathbf{y}^2\| + 2t(\mathbf{x}, \mathbf{y}) + \|\mathbf{x}\|^2$$

Thus, the discriminant must be less than zero (because the whole polynomial is positive, so the discriminant [the opposite of the x^0 term of the factored form of the polynomial] must be less than zero so the polynomial doesn't get dragged down to negative values):

$$(\mathbf{x}, \mathbf{y})^2 - \|\mathbf{x}\|^2 \|\mathbf{y}\|^2 \le 0$$

Taking square roots of both sides proves the desired inequality.

- Recall that if $A^* = A$, then all eigenvalues are real and all eigenvectors of distinct eigenvalues are orthogonal to each other.
- Normal (matrix): A matrix N such that $N^*N = NN^*$.
 - Examples: Diagonal, self-adjoint, and unitary operators are all normal.
- Any normal operator in a complex vector spae has an orthonormal set of eigenvectors, e.g., $N = UDU^*$.
 - Proof: N is upper triangular wrt. some basis (because all matrices are). WTS any normal upper triangular matrix is diagonal. Done by induction on the dimension of N from n = 2.
 - Assume the claim for every $(n-1) \times (n-1)$ normal upper triangular matrix.
 - Let

$$N = \begin{pmatrix} a_{11} & a_{12} & \cdots & a_{1n} \\ 0 & & & \\ 0 & & & \\ \vdots & & N_1 & \\ 0 & & & \end{pmatrix}$$

(we know every normal matrix can be written in this upper triangular form)

- Then just compute NN^* and N^*N . Knowing they have to be equal, we have that $a_{12} = \cdots = a_{1n} = 0$.
- We can also prove from the above (block diagonal multiplication) that N_1 is normal. Thus, it's diagonal, too. Therefore, the whole thing is diagonal.
- N is normal if and only if $||N\mathbf{x}|| = ||N^*\mathbf{x}||$.
 - Proof: $(N\mathbf{x}, N\mathbf{y}) = (N^*N\mathbf{x}, \mathbf{y}) = (NN^*\mathbf{x}, \mathbf{y}) = (N^*\mathbf{x}, N^*\mathbf{y})$. This is equivalent to the desired condition.
- If A is nonnegative and $(A\mathbf{e}_k, \mathbf{e}_k) = a_{kk}$, then

$$\sum_{i,j=1}^{n} a_{ij} \mathbf{x}_i \mathbf{x}_j$$

- Positive definite (matrix): An $n \times n$ self-adjoint matrix such that $(A\mathbf{x}, \mathbf{x}) > 0$ for all $\mathbf{x} \in X$.
- Let $A: X \to Y$, dim $X = \dim Y$. Then AA^* is positive semidefinite. And there exists a unique square root $R = \sqrt{A^*A}$.
 - Proof: $(A^*A\mathbf{x}, \mathbf{x}) = (A\mathbf{x}, A\mathbf{x}) = ||A\mathbf{x}||^2 \ge 0.$
- Modulus (of A): The matrix $|A| = \sqrt{A^*A}$.
- Check $||A|\mathbf{x}|| = ||A\mathbf{x}||$.

$$|||A|\mathbf{x}||^2 = (|A|\mathbf{x}, |A|\mathbf{x}) = (|A|^*|A|\mathbf{x}, \mathbf{x}) = (A^*A\mathbf{x}, \mathbf{x}) = (A\mathbf{x}, A\mathbf{x}) = ||A\mathbf{x}||^2$$

- Let $A: X \to X$ be a linear operator. Then A = U|A| where U is unitary.
- Look at singular matrices.
- Recall that if $A: X \to Y$, we have that A^*A is semidefinite, positive, and self adjoint.
 - Thus, there exists a unique matrix $R = \sqrt{A^*A} \ge 0$, which we define to be $|A| = \sqrt{A^*A}$.
 - Polar form of a matrix:

10/15:

$$A = U|A|$$

- This may not be unique!
- Proof: Suppose $A\mathbf{x} = U(|A|\mathbf{x})$. $A\mathbf{x} \in \text{range } A$, and $|A|\mathbf{x} \in \text{range}(|A|)$. $\mathbf{x} \in \text{range}(|A|)$ implies that there exists $\mathbf{v} \in X$ such that $x = |A|\mathbf{v}$.
- Define $U\mathbf{x} = A\mathbf{x}$. U is a well-defined linear map.
- $\|U_0 \mathbf{x}\| = \|A\mathbf{x}\| = \||A|\mathbf{v}\| = \|\mathbf{x}\|.$
- U is an isometry.
- range $|A| \to X$.
- Use ker $A = \ker |A| = (\operatorname{range} A)^{\perp}$ to extend U_0 to $U: U = U_0 + U_1$.
- Singular values (of a matrix): The eigenvalues of |A|.
 - So if $\lambda_1, \ldots, \lambda_n$ are the eigenvalues of A^*A , the singular values of A are $\sqrt{\lambda_1}, \ldots, \sqrt{\lambda_n}$.
- Let $A: X \to Y$ be a linear map.
 - Let $\sigma_1, \ldots, \sigma_n$ be the signular values of A. Then $\sigma_1, \ldots, \sigma_n > 0$.
 - Additionally, if $\mathbf{v}_1, \dots, \mathbf{v}_n$ is an orthonormal basis of eigenvectors of A^*A , then the list of n vectors $\mathbf{w}_1, \dots, \mathbf{w}_n$ defined by $\mathbf{w}_k = 1/\sigma_k A \mathbf{v}_k$ for each $k = 1, \dots, n$ is orthonormal.

■ Proof:

$$(\mathbf{w}_k, \mathbf{w}_j) = \frac{1}{\sigma_k \sigma_k} (A\mathbf{v}_k, A\mathbf{v}_j) = \frac{1}{\sigma_k \sigma_j} = \frac{1}{\sigma_k \sigma_j} (A^* A\mathbf{v}_k, \mathbf{v}_j) = \frac{\sigma_k^2}{\sigma_k \sigma_j} (\mathbf{v}_k, \mathbf{v}_j) = 0$$

and

$$\|\mathbf{w}_k\| = \frac{1}{\sigma_k} \|A\mathbf{v}_k\| = \frac{1}{\sigma_k} \||A|\mathbf{v}_k\| = 1$$

- Schmidt decomposition of A:

$$A\mathbf{x} = \sum_{k=0}^{r} \sigma_k(\mathbf{x}, \mathbf{v}_k) \mathbf{w}_k$$

■ This is because $\mathbf{x} = \sum (\mathbf{x}, \mathbf{v}_k) \mathbf{v}_k$, so by the above,

$$A\mathbf{x} = \sum_{k=0}^{n} (\mathbf{x}, \mathbf{v}_{k}) A\mathbf{v}_{k} = \sum_{k=0}^{r} \sigma_{k}(\mathbf{x}, \mathbf{v}_{k}) \mathbf{w}_{k}$$

- Operator norm: $||A|| = \max\{||A\mathbf{x}|| : ||\mathbf{x}|| \le 1\}.$
- Properties of the operator norm:
 - $\|A\mathbf{x}\| \le \|A\| \|\mathbf{x}\|.$
 - $\|\alpha A\| = |\alpha| \|A\|.$
 - $\|A + B\| \le \|A\| + \|B\|.$
 - $\|A\| \ge 0.$
 - $\|A\| = 0 \text{ iff } A = 0.$
- Frobenius norm: The norm $||A||_2^2 = \operatorname{tr}(A^*A)$.
- The operator norm is always less than or equal to the Frobenius norm.
- If $A: \mathbb{F}^n \to \mathbb{F}^n$, then $A = W \Sigma V^*$ where σ is a diagonal matrix of nonzero singular values.
- The operator norm of A is the largest of the singular values.
- An orthogonal matrix can be decomposed to a block-diagonal matrix of rotations.
- 10/18: Soug tests what he teaches and doesn't give super tricky questions.
 - Structure of orthogonal matrices.
 - Orthogonal (matrix): A unitary matrix U with all elements real and $|\det U| = 1$.
 - Theorem: Let U be an orthogonal operator on \mathbb{R}^n such that $\det U = 1$. Then there exists an orthonormal basis $\mathbf{v}_1, \dots, \mathbf{v}_n$ such that with respect to this basis,

$$U = \begin{pmatrix} R_{\phi_1} & \mathbf{0} \\ & \ddots & \\ & & \mathbb{R}_{\phi_k} \\ \mathbf{0} & & I_{n-2k} \end{pmatrix}$$

where each R_{ϕ_i} is a 2 × 2 rotation matrix.

- If you are in \mathbb{R}^7 for example, you would be able to express U as a composition of at most 3 rotation maps and the identity map.
- Each rotation map acts on two orthonormal vectors.
- Proof: $P(\lambda)$ is the *n*-degree characteristic polynomial $\det(U \lambda I) = 0$. The eigenvalues are the roots of it.

- $-p(\lambda)=0$ if and only if $p(\bar{\lambda})=0$.
 - $\lambda \in \mathbb{C}$ is an eigenvalue with eigenvector $\mathbf{u} \neq 0$ iff $U\mathbf{u} = \lambda \mathbf{u}$ and $U\bar{\mathbf{u}} = \bar{\lambda}\bar{\mathbf{u}}$.
- Recall that U unitary implies $|\lambda| = 1$.
 - Proof^[1]: $||U\mathbf{x}|| = ||\mathbf{x}||$ and $U\mathbf{x} = \lambda \mathbf{x}$. Thus,

$$||U\mathbf{x}|| = ||\lambda\mathbf{x}|| = |\lambda|||\mathbf{x}|| = ||\mathbf{x}||$$

and since $\mathbf{x} \neq 0$, we can divide by $\|\mathbf{x}\|$, so $|\lambda| = 1$.

- $\operatorname{Let} \mathbf{u} = \operatorname{Re} \mathbf{u} + \operatorname{Im} \mathbf{u}.$
- It follows that we may define

$$\mathbf{x} = \operatorname{Re} \mathbf{u} = \frac{\mathbf{u} + \bar{\mathbf{u}}}{2}$$
 $\mathbf{y} = \operatorname{Im} \mathbf{u} = \frac{\mathbf{u} - \bar{\mathbf{u}}}{2}$

- Thus, $\mathbf{u} = \mathbf{x} + i\mathbf{y}$ and $\bar{\mathbf{u}} = \mathbf{x} i\mathbf{y}$.
- Since $U\mathbf{x} = \frac{U\mathbf{u} + U\bar{\mathbf{u}}}{2} = \frac{\lambda \mathbf{u} + \bar{\lambda}\bar{\mathbf{u}}}{2}$, $U\mathbf{y} = \text{Im}(\lambda \mathbf{u}) = \text{Re}(\lambda \mathbf{u})$.
- Since $|\lambda| = 1$, $\lambda = e^{i\alpha}$ and $\bar{\lambda} = e^{-i\alpha}$.
- It follows that $U\mathbf{x} = (\cos \alpha)\mathbf{x} (\sin \alpha)\mathbf{y}$ and $U\mathbf{y} = (\cos \alpha)\mathbf{y} + (\sin \alpha)\mathbf{x}$.
- Thus, since $U\mathbf{x} = \operatorname{Re} \lambda \mathbf{u}$, we have that

$$\lambda \mathbf{u} = (\cos \alpha + i \sin \alpha)(\mathbf{x} + i\mathbf{y})$$

= $(\cos \alpha)\mathbf{x} - (\sin \alpha)\mathbf{y} + i[(\cos \alpha)\mathbf{y} + (\sin \alpha)\mathbf{x}]$

- If E_{λ} is a 2 dimensional space spanned by **x** and **y** and invariant by U. Thus, any block of the desired matrix leaves its desired sub-block invariant.
- We also know that the eigenvectors of a unitary matrix corresponding to different eigenvalues are orthogonal.
- Thus, $\|\mathbf{x}\| = \|\mathbf{y}\| = \sqrt{2}/2\|\mathbf{u}\|, \mathbf{x} \perp \mathbf{y}.$
- Let \mathbf{x}, \mathbf{y} complete the theorem to form a basis of \mathbb{R}^n .
- It will follow that

$$U = \begin{pmatrix} R_{\alpha} & \mathbf{0} \\ \mathbf{0} & U_{1} \end{pmatrix}$$

where U_1 is orthogonal, and we may repeat the process.

6.2 Chapter 6: Structure of Operators on Inner Product Spaces

From Treil (2017).

10/24:

- Theorem 6.1.1: Let $A: X \to X$ be an operator acting in a complex inner product space. Then there exists an orthonormal basis $\mathbf{u}_1, \dots, \mathbf{u}_n$ of X such that the matrix of A in this basis is upper triangular. In other words, any $n \times n$ matrix A can be represented as $A = UTU^*$, where U is unitary and T is upper-triangular.
- Theorem 6.1.2: Let $A: X \to X$ be an operator acting on a real inner product space. Suppose that all eigenvalues of A are real. Then there exists an orthonormal basis $\mathbf{u}_1, \ldots, \mathbf{u}_n$ in X such that the matrix of A in this basis is upper triangular. In other words, any real $n \times n$ matrix A with all real eigenvalues can be represented as $T = UTU^* = UTU^T$, where U is orthogonal and T is a real upper-triangular matrix.

 $^{^{1}}$ This would be a good exam question.

- Theorem 6.2.1: Let $A = A^*$ be a self-adjoint operator in an inner product space X (the space can be complex or real). Then all eigenvalues of A are real and there exists an orthonormal basis of eigenvectors of A in X.
 - Equivalently (see Theorem 6.2.2), A can be represented as $A = UDU^*$ where U is a unitary matrix and D is a diagonal matrix with real entries. Moreover, if A is real, U can be chosen to be real, i.e., orthogonal.
- Proposition 6.2.3: Let $A = A^*$ be a self-adjoint operator and let $\lambda, \mathbf{u}, \mu, \mathbf{v}$ be such that $A\mathbf{u} = \lambda \mathbf{u}$ and $A\mathbf{v} = \mu \mathbf{v}$. Then if $\lambda \neq \mu, \mathbf{u} \perp \mathbf{v}$.
- Since complex multiplication is commutative,

$$D^*D = DD^*$$

for every diagonal matrix D.

- It follows that $A^*A = AA^*$ if the matrix of A in some orthonormal basis is diagonal.
- Theorem 6.2.4: Any normal operator N in a complex vector space has an orthonormal basis of eigenvectors.
 - Equivalently, any matrix N satisfying $N^*N=NN^*$ can be represented as $N=UDU^*$ where U is unitary and D is diagonal.
- Proposition 6.2.5: An operator $N: X \to X$ is normal iff

$$||N\mathbf{x}|| = ||N^*\mathbf{x}||$$

for all $\mathbf{x} \in X$.

- Hermitian square (of A): The matrix A^*A .
- Modulus (of A): The unique positive semidefinite square root $\sqrt{A^*A}$.
- Proposition 6.3.3: For a linear operator $A: X \to Y$,

$$|||A|\mathbf{x}|| = ||A\mathbf{x}||$$

- Corollary 6.3.4: $\ker A = \ker |A|$.
- Theorem 6.3.5: Let $A: X \to X$ be an operator (square matrix). Then A can be represented as

$$A = U|A|$$

where U is a unitary operator.

- Singular value (of A): An eigenvalue of |A|.
 - A positive square root of an operator of A^*A .
- Proposition 6.3.6: Let $\sigma_1, \ldots, \sigma_n$ be the singular values of A, ordered such that $\sigma_1, \ldots, \sigma_r$ are the nonzero singular values, and let $\mathbf{v}_1, \ldots, \mathbf{v}_n$ be an orthonormal basis of eigenvectors of A^*A . Then the system

$$\mathbf{w}_k = \frac{1}{\sigma_k} A \mathbf{v}_k$$

for k = 1, ..., r is orthonormal.

• Schmidt decomposition (of A): The decompositions

$$A = \sum_{k=1}^{r} \sigma_k \mathbf{w}_k \mathbf{v}_k^*$$

and

10/25:

$$A\mathbf{x} = \sum_{k=1}^{r} \sigma_k(\mathbf{x}, \mathbf{v}_k) \mathbf{w}_k$$

- Note that these can be verified by plugging $\mathbf{x} = \mathbf{v}_j$ for each $j = 1, \dots, n$ into the latter equation.
- Lemma 6.3.7: A can be represented as the Schmidt decomposition

$$A = \sum_{k=1}^{r} \sigma_k \mathbf{w}_k \mathbf{v}_k^*$$

where $\sigma_k > 0$ for any orthonormal systems $\mathbf{v}_1, \dots, \mathbf{v}_r$ and $\mathbf{w}_1, \dots, \mathbf{w}_r$.

• Corollary 6.3.8: Let $A = \sum_{k=1}^{r} \sigma_k \mathbf{w}_k \mathbf{v}_k^*$ be a Schmidt decomposition of A. Then

$$A^* = \sum_{k=1}^r \sigma_k \mathbf{v}_k \mathbf{w}_k^*$$

is a Schmidt decomposition of A^* .

• Reduced singular value decomposition (of A): The decomposition

$$A = \tilde{W}\tilde{\Sigma}\tilde{V}^*$$

where $A: \mathbb{F}^n \to \mathbb{F}^m$ has the Schmidt decomposition $A = \sum_{k=1}^r \sigma_k \mathbf{w}_k \mathbf{v}_k^*$, $\tilde{\Sigma} = \text{diag}\{\sigma_1, \dots, \sigma_r\}$, and \tilde{V}, \tilde{W} are matrices with columns $\mathbf{v}_1, \dots, \mathbf{v}_r$ and $\mathbf{w}_1, \dots, \mathbf{w}_r$, respectively. Also known as compact singular value decomposition.

- Note that \tilde{V} is an $n \times r$ matrix, $\tilde{\Sigma}$ is an $r \times r$ matrix, and \tilde{W} is an $m \times r$ matrix.
- Since $\mathbf{v}_1, \dots, \mathbf{v}_r$ and $\mathbf{w}_1, \dots, \mathbf{w}_r$ are orthonormal, \tilde{V}, \tilde{W} are isometries.
- Note that $r = \operatorname{rank} A$ (see Problem 6.3.1).
 - It follows that if A is invertible, then m=n=r, so \tilde{V},\tilde{W} are unitary and $\tilde{\Sigma}$ is an invertible diagonal matrix.
- However, A need not be invertible for us to get a representation similar to $A = \tilde{W} \tilde{\Sigma} \tilde{V}^*$.
 - Complete $\mathbf{v}_1, \dots, \mathbf{v}_r$ and $\mathbf{w}_1, \dots, \mathbf{w}_r$ to bases of \mathbb{F}^n and \mathbb{F}^m , respectively.
 - Then we get the following.
- Singular value decomposition (of A): The decomposition

$$A = W\Sigma V^*$$

where $V \in M_{n \times n}^{\mathbb{F}}$ and $W \in M_{m \times m}^{\mathbb{F}}$ are unitary matrices with columns $\mathbf{v}_1, \dots, \mathbf{v}_n$ and $\mathbf{w}_1, \dots, \mathbf{w}_n$, respectively, and $\Sigma \in M_{m \times n}^{\mathbb{R}^+}$ is a "diagonal" matrix such that

$$\Sigma_{j,k} = \begin{cases} \sigma_k & j = k \le r \\ 0 & \text{otherwise} \end{cases}$$

• Notice that if $A = W\Sigma V^*$, then

$$A^*A = (W\Sigma V^*)^*(W\Sigma V^*) = V\Sigma^*W^*W\Sigma V^* = V\Sigma^2V^*$$

proving that the singular values of A, squared, are the eigenvalues of A^*A .

- If A is invertible, the reduced SVD is the matrix form of the Schmidt decomposition is the SVD.
- If $A = W\Sigma V^*$ is $n \times n$, then

$$A = (\underbrace{WV^*}_{U})(\underbrace{V\Sigma V^*}_{|A|})$$

is a polar decomposition of A.

- Consider the unit ball $B = \{ \mathbf{x} \in \mathbb{R}^n : ||\mathbf{x}|| \le 1 \}.$
 - We want to describe A(B), i.e., the image of the unit ball under A.
 - Let $\mathbf{x} = (x_1, \dots, x_n)^T$ and let $\mathbf{y} = (y_1, \dots, y_n)^T$. If $A = \text{diag}\{\sigma_1, \dots, \sigma_n\}$, we have $\mathbf{y} \in A(B)$ iff $\mathbf{y} = A\mathbf{x}$ where $\mathbf{x} \in B$ iff

$$\sum_{k=1}^{n} \frac{y_k^2}{\sigma_k^2} = \sum_{k=1}^{n} x_k^2 = \|\mathbf{x}\|^2 \le 1$$

- Thus, A(B) is an ellipsoid with half-axes $\sigma_1, \ldots, \sigma_n$.
- In the more general case, if $A = W\Sigma V^*$, then since V^* is unitary, $V^*(B) = B$. $\Sigma V^*(B) = \Sigma(B)$ is thus by the above an ellipsoid in range Σ with half-axes $\sigma_1, \ldots, \sigma_r$. Thus, since isometries don't change geometry, $W(\Sigma(B))$ is also an ellipsoid with the same half-axes, but in range A.
- Conclusion: The image A(B) of the closed unit ball B is an ellipsoid in range A with half-axes $\sigma_1, \ldots, \sigma_r$, where r is the number of nonzero singular values, i.e., the rank of A.
- Finding the maximum of $||A\mathbf{x}||$ for $\mathbf{x} \in B$.
 - For a diagonal matrix Σ with nonnegative entries, the maximum is clearly the maximal diagonal entry: In this case if s_1 is the maximal diagonal entry, then since

$$\Sigma \mathbf{x} = \sum_{k=1}^{r} s_k x_k \mathbf{e}_k$$

we have that

$$||A\mathbf{x}||^2 = \sum_{k=1}^r s_k^2 |x_k|^2 \le s_1^2 \sum_{k=1}^r |x_k|^2 = s_1^2 \cdot ||\mathbf{x}||^2$$

- We get the following by a similar logic to before.
- Conclusion: The maximum of $||A\mathbf{x}||$ on the unit ball B is the maximal singular value of A.
- Operator norm (of A): The following quantity. Denoted by ||A||. Given by

$$||A|| = \max\{||A\mathbf{x}|| : \mathbf{x} \in X, ||\mathbf{x}|| \le 1\}$$

- $\|A\|$ clearly satisfies the four properties of a norm.
- Additionally,

$$||A\mathbf{x}|| < ||A|| \cdot ||\mathbf{x}||$$

– Alternate definition: The operator norm ||A|| is the smallest number $C \ge 0$ such that $||A\mathbf{x}|| \le C||\mathbf{x}||$.

• Frobenius norm: The following norm. Also known as Hilbert-Schmidt norm. Denoted by $\|A\|_2$. Given by

$$||A||_2^2 = \operatorname{tr}(A^*A)$$

- If we let s_1, \ldots, s_n be the singular values of A and let s_1 be the largest value, then we have

$$||A||^2 = s_1^2 \le \sum_{k=1}^n s_k^2 = \operatorname{tr}(A^*A) = ||A||_2^2$$

- Conclusion: The operator norm of a matrix cannot be more than its Frobenius norm.
- Suppose we want to solve $A\mathbf{x} = \mathbf{b}$ where A is invertible, but there is some (experimental) error $\Delta \mathbf{b}$ in \mathbf{b} . Then we are really solving for an approximate solution $\mathbf{x} + \Delta \mathbf{x}$ to the equation

$$A(\mathbf{x} + \Delta \mathbf{x}) = \mathbf{b} + \Delta \mathbf{b}$$

- It follows since A is invertible that $\mathbf{x} = A^{-1}\mathbf{b}$ and $\Delta \mathbf{x} = A^{-1}\Delta \mathbf{b}$.
- To estimate the relative error $\|\Delta \mathbf{x}\|/\|\mathbf{x}\|$ in the solution in comparison with the relative error $\|\Delta \mathbf{b}\|/\|\mathbf{b}\|$ in the data, use

$$\frac{\|\Delta \mathbf{x}\|}{\|\mathbf{x}\|} = \frac{\|A^{-1}\Delta \mathbf{b}\|}{\|\mathbf{b}\|} \frac{\|A\mathbf{x}\|}{\|\mathbf{x}\|} \le \frac{\|A^{-1}\| \cdot \|\Delta \mathbf{b}\|}{\|\mathbf{b}\|} \frac{\|A\| \cdot \|\mathbf{x}\|}{\|\mathbf{x}\|} = \|A^{-1}\| \cdot \|A\| \cdot \frac{\|\Delta \mathbf{b}\|}{\|\mathbf{b}\|}$$

• Condition number (of A): The following quantity. Given by

$$||A|| \cdot ||A^{-1}||$$

- If s_1 is the largest singular value of A and s_n is the smallest, then

$$||A|| \cdot ||A^{-1}|| = s_1 \cdot \frac{1}{s_n} = \frac{s_1}{s_n}$$

- Well-conditioned (matrix): A matrix the condition number of which is not "too big."
- Ill-conditioned (matrix): A matrix that is not well-conditioned.
- Theorem 6.5.1: Let U be an orthogonal operator on \mathbb{R}^n and let $\det U = 1$. Then there exists an orthonormal basis $\mathbf{v}_1, \dots, \mathbf{v}_n$ such that the matrix of U in this basis has the block diagonal form

$$\begin{pmatrix} R_{\varphi_1} & & 0 \\ & \ddots & \\ & & R_{\varphi_k} \\ 0 & & I_{n-2k} \end{pmatrix}$$

where each R_{φ_i} is a two-dimensional rotation

$$R_{\varphi_j} = \begin{pmatrix} \cos \varphi_j & -\sin \varphi_j \\ \sin \varphi_j & \cos \varphi_j \end{pmatrix}$$

and I_{n-2k} represents the $(n-2k) \times (n-2k)$ identity matrix.

- Alternate interpretation: Any rotation in \mathbb{R}^n can be represented as a composition of at most n/2 commuting planar rotations.

• Theorem 6.5.2: Let U be an orthogonal operator on \mathbb{R}^n and let $\det U = -1$. Then there exists an orthonormal basis $\mathbf{v}_1, \dots, \mathbf{v}_n$ such that the matrix of U in this basis has block diagonal form

$$\begin{pmatrix} R_{\varphi_1} & & & 0 \\ & \ddots & & & \\ & & R_{\varphi_k} & & \\ & & & I_r & \\ 0 & & & -1 \end{pmatrix}$$

where r = n - 2k - 1 and each R_{φ_i} is a two-dimensional rotation

$$R_{\varphi_j} = \begin{pmatrix} \cos \varphi_j & -\sin \varphi_j \\ \sin \varphi_j & \cos \varphi_j \end{pmatrix}$$

- Corollary: An orthogonal 2×2 matrix U with determinant -1 is always a reflection.
- Theorem 6.5.3: Any rotation U (i.e., any orthogonal transformation U with $\det U = 1$) can be represented as a product of at most n(n-1)/2 elementary rotations.
- Consider the following orthonormal bases of \mathbb{R}^2 .



Figure 6.1: Orientation in \mathbb{R}^2 .

- Notice that a rotation will get you from the standard basis (a) to basis (b), but not from the standard basis (a) to basis (c).
- This is the motivation for defining orientation.
- More formally, we know that there is a unique linear transformation U such that $U\mathbf{e}_k = \mathbf{v}_k$ for each k = 1, 2. In particular, the matrix of U with respect to the standard basis is orthogonal with columns $\mathbf{v}_1, \mathbf{v}_2$.
- By Theorems 6.5.1 and 6.5.2, if det U = 1, then U is a rotation, and if det U = -1, then U is not a rotation.
- Similarly oriented (bases \mathcal{A}, \mathcal{B}): Two bases \mathcal{A}, \mathcal{B} of a real vector space such that the change of coordinates matrix $[I]_{\mathcal{B}\mathcal{A}}$ has a positive determinant.
- **Differently oriented** (bases \mathcal{A}, \mathcal{B}): Two bases \mathcal{A}, \mathcal{B} of a real vector space that are not similarly oriented (i.e., $[I]_{\mathcal{B}\mathcal{A}}$ has a negative determinant).
- We usually let the standard basis of \mathbb{R}^n have a **positive orientation**.
 - In an abstract vector space, we need only fix a basis and declare its orientation to be positive.
- Continuously transformable (bases \mathcal{A}, \mathcal{B}): Two bases \mathcal{A}, \mathcal{B} such that $\mathcal{A} = \{\mathbf{a}_1, \dots, \mathbf{a}_n\}$ can be continuously transformed to a basis $\mathcal{B} = \{\mathbf{b}_1, \dots, \mathbf{b}_n\}$. In particular, there exists a **continuous** family of bases $\mathcal{V}(t) = \{\mathbf{v}_1(t), \dots, \mathbf{v}_n(t)\}$, $t \in [a, b]$, such that

$$\mathbf{v}_k(a) = \mathbf{a}_k \qquad \qquad \mathbf{v}_k(b) = \mathbf{b}_k$$

for each $k = 1, \ldots, n$.

- Continuous family of bases: A family of bases $\mathcal{V}(t) = \{\mathbf{v}_1(t), \dots, \mathbf{v}_n(t)\}, t \in [a, b]$, such that the vector-functions $\mathbf{v}_k(t)$ are continuous (their coordinates in some bases are continuous functions) and the system $\mathbf{v}_1(t), \dots, \mathbf{v}_n(t)$ is a basis for all $t \in [a, b]$.
- Theorem 6.6.1: Two bases $\mathcal{A} = \{\mathbf{a}_1, \dots, \mathbf{a}_n\}$ and $\mathcal{B} = \{\mathbf{b}_1, \dots, \mathbf{b}_n\}$ have the same orientation if and only if one of the bases can be continuously transformed to the other.

Bilinear and Quadratic Forms

7.1 Notes

• Bilinear form: A function $L: \mathbb{R}^n \times \mathbb{R}^n \to \mathbb{R}$ such that 10/18:

$$L(\alpha \mathbf{x}_1 + \beta \mathbf{x}_2, \mathbf{y}) = \alpha L(\mathbf{x}_1, \mathbf{y}) + \beta L(\mathbf{x}_2, \mathbf{y}) \qquad L(\mathbf{x}, \alpha \mathbf{y}_1 + \beta \mathbf{y}_2) = \alpha L(\mathbf{x}, \mathbf{y}_1) + \beta L(\mathbf{x}, \mathbf{y}_2)$$

$$-L(\mathbf{x}, \mathbf{y}) = (A\mathbf{x}, \mathbf{y}).$$

• Quadratic form: A bilinear form $L(\mathbf{x}, \mathbf{x})$.

 $-(\mathbf{x},\mathbf{x})$ is a polynomial of degree 2 in $\mathbf{x}_1,\ldots,\mathbf{x}_n$:

$$L(\lambda \mathbf{x}, \lambda \mathbf{x}) = (\lambda \mathbf{x}, \lambda \mathbf{x}) = \lambda^2(\mathbf{x}, \mathbf{x})$$

• We have that

$$(A\mathbf{x}, \mathbf{x}) = (A\lambda\mathbf{x}, \lambda\mathbf{x}) = \lambda^2(A\mathbf{x}, \mathbf{x}) = \sum_{j,i=1}^n \alpha_{j,i}\mathbf{x}_i\mathbf{x}_j$$

• The general form of a quadratic form:

- Can any quadratic form on \mathbb{R}^n be written as $(A\mathbf{x}, \mathbf{x})$?

10/20: • Bilinear forms are linear in each argument when keeping the other fixed.

• Quadratic forms $Q(\mathbf{x}) = L(\mathbf{x}, \mathbf{x})$ are quadratic polynomials in the coordinates of x.

- In particular, $Q(\lambda \mathbf{x}) = |\lambda|^2 Q(\mathbf{x})$.

• If Q quadratic is real, then $Q(\mathbf{x}) = (A\mathbf{x}, \mathbf{x})$ where A is some square matrix.

- If $\mathbf{e}_1, \dots, \mathbf{e}_n$ is an orthonormal basis of \mathbb{R}^n , then there exists a unique $A = A^*$ such that $(A)_{ij} =$

- Keeping $\mathbf{x} = \sum_{i=1}^{n} \mathbf{x}_i, \mathbf{e}_i$ foxed, we have

$$\begin{aligned} Q(\mathbf{x}) &= L(\mathbf{x}, \mathbf{x}) \\ &= L(\sum_{i=1}^{n} \mathbf{x}_{i} \mathbf{e}_{i}, \sum_{i=1}^{n} \mathbf{x}_{j} \mathbf{e}_{j}) \\ &= \sum_{i=1}^{n} \mathbf{x}_{i} L(\mathbf{e}_{i}, \sum_{i=1}^{n} \mathbf{x}_{j} \mathbf{e}_{j}) \\ &= \sum_{i,j=1}^{n} \mathbf{x}_{i} \mathbf{x}_{j} \underbrace{L(\mathbf{e}_{i}, \mathbf{e}_{j})}_{A_{ij}} \end{aligned}$$

• We have that

$$(A\mathbf{x}, \mathbf{x}) = (UDU^{-1}\mathbf{x}, \mathbf{x})$$

$$= (DU^{-1}\mathbf{x}, U^{-1}\mathbf{x})$$

$$= \sum_{i=1}^{n} \lambda_{i} (\underbrace{U^{-1}\mathbf{x}}_{\mathbf{y}_{i}})_{i} (\underbrace{U^{-1}\mathbf{x}}_{\mathbf{y}_{i}})_{i}$$

- Can we characterize the set $\{\mathbf{x}: (A\mathbf{x}, \mathbf{x}) = 1\}$?
 - Note that this set is equivalent to $\{\mathbf{y}:(D\mathbf{y},\mathbf{y})=1\}$ by teh above. This set is a rotation of the previous one. Ellipse?
- Positive quadratic form:
 - Q is positive definite if $Q(\mathbf{x}) > \mathbf{0}$ for all $\mathbf{x} \neq \mathbf{0}$ and Q is positive semidefinite if $Q(\mathbf{x}) \geq \mathbf{0}$ for all $\mathbf{x} \neq \mathbf{0}$.
 - Take a self-adjoint matrix $A = A^*$. It is positive definite if $Q(\mathbf{x}) = (A\mathbf{x}, \mathbf{x})$ is positive definite.
- Theorem: If $A = A^*$, then
 - 1. A is positive definite if and only if all eigenvalues of A are positive.
 - 2. A is positive semidefinite if and only if all eigenvalues of A are nonnegative.
 - 3. A is negative semidefinite if and only if all eigenvalues of A are nonpositive.
 - 4. A is negative definite if and only if all eigenvalues of A are negative.
 - 5. A is indefinite if and only if the eigenvalues of A have positive and negative values.
- Theorem: $A = A^*$ is positive definite iff det $A_k > 0$ for all k = 1, ..., n where A_k is the upper left $k \times k$ submatrix.
- Minimax representation of eigenvalues of a self-adjoint A.
 - Let E be a subspace of X where dim $X < \infty$. We define $\operatorname{codim}(E) = \dim E^{\perp}$.
 - Thus, $\dim E + \operatorname{codim} E = \dim X$.
 - Theorem: Let $A=A^*,\ \lambda_1\geq\cdots\geq\lambda_n$ eigenvalues of A. Then

$$\lambda_k = \max_{\substack{\text{E subspace} \\ \dim E = k}} \min_{\substack{\mathbf{x} \in E \\ \|\mathbf{x}\| = 1}} (A\mathbf{x}, \mathbf{x}) = \min_{\substack{\text{F subspace} \\ \operatorname{codim} F = k - 1}} \max_{\substack{\mathbf{x} \in F \\ \|\mathbf{x}\| = 1}} (A\mathbf{x}, \mathbf{x})$$

- Proof: A diagonal equals $(\lambda_1, \ldots, \lambda_n)$.
- An orthonormal basis of X such that dim E = k, codim F = k 1, dim F = n k + 1.
- There exists an $\mathbf{x}_0 \neq \mathbf{0}$ such that $\mathbf{x}_0 \in E \cap F$.
- Note that if $B = B^*$, then the max and min of $(B\mathbf{x}, \mathbf{x})$ over the unit sphere is the maximal and minimal eigenvalue of B.
- Thus,

$$\min_{\substack{\mathbf{x} \in E \\ \|\mathbf{x}\|=1}} (A\mathbf{x}, \mathbf{x}) \le (A\mathbf{x}_0, \mathbf{x}_0) \le \max_{\substack{\mathbf{x} \in F \\ \|\mathbf{x}\|=1}} (A\mathbf{x}, \mathbf{x})$$

- This is true for any E, F subspaces. dim E = k, codim F = k 1, $E_0 = \operatorname{span}(\mathbf{e}_1, \dots, \mathbf{e}_k)$ and $F_0 = \operatorname{span}(\mathbf{e}_k, \dots, \mathbf{e}_n)$.
- Thus,

$$\min_{\substack{E_0\\ \|\mathbf{x}\|=1}} (A\mathbf{x}, \mathbf{x}) = \lambda_k = \max_{\substack{F_0\\ \|\mathbf{x}\|=1}} (A\mathbf{x}, \mathbf{x})$$

■ Additionally,

$$\lambda_{k_1} \leq \max_{\dim E=k} \min_{\mathbf{x}} (A\mathbf{x}, \mathbf{x}) \leq \min_{\substack{F \ \text{codim } F=k-1}} \max_{\mathbf{x}} (A\mathbf{x}, \mathbf{x}) \leq \lambda_k$$

- Corollary: Let $A = A^* = (a_{jk})_{1 \leq j,k \leq n}$ with eigenvalues $\lambda_1, \ldots, \lambda_n$ listed in decreasing order. Let $\tilde{A} = (a_{j,k})_{1 \leq j,k \leq n-1}$ with eigenvalues μ_1, \ldots, μ_{n-1} listed in decreasing order. Then $\lambda_1 \geq \mu_1 \geq \lambda_2 \geq \mu_2 \geq \cdots \geq \mu_{n-1} \geq \lambda_n$.
 - Consider $(A\mathbf{x}, \mathbf{x})$ on $\{\mathbf{e}_1, \dots, \mathbf{e}_n\}$, but then restrict yourself to $\mathbf{x} \in \mathbb{R}^{n-1}$ on $\{\mathbf{e}_1, \dots, \mathbf{e}_{n-1}\}$.

7.2 Chapter 7: Bilinear and Quadratic Forms

From Treil (2017).

10/25:

- Bilinear form (on \mathbb{R}^n): A function $L(\mathbf{x}, \mathbf{y})$ of two arguments $\mathbf{x}, \mathbf{y} \in \mathbb{R}^n$ that is linear in each argument.
 - Linearity in each argument:

$$L(\alpha \mathbf{x}_1 + \beta \mathbf{x}_2, \mathbf{y}) = \alpha L(\mathbf{x}_1, \mathbf{y}) + \beta L(\mathbf{x}_2, \mathbf{y})$$
 $L(\mathbf{x}, \alpha \mathbf{y}_1 + \beta \mathbf{y}_2) = \alpha L(\mathbf{x}, \mathbf{y}_1) + \beta L(\mathbf{x}, \mathbf{y}_2)$

• If $\mathbf{x} = (x_1, \dots, x_n)^T$ and $\mathbf{y} = (y_1, \dots, y_n)^T$, then

$$L(\mathbf{x}, \mathbf{y}) = \sum_{j,k=1}^{n} a_{j,k} x_k y_j$$
$$= (A\mathbf{x}, \mathbf{y})$$
$$= \mathbf{y}^T A\mathbf{x}$$

where

$$A = \begin{pmatrix} a_{1,1} & \cdots & a_{1,n} \\ \vdots & & \vdots \\ a_{n,1} & \cdots & a_{n,n} \end{pmatrix}$$

- -A is uniquely determined by L.
- Quadratic form (on \mathbb{R}^n): The diagonal of a bilinear form L, i.e., a bilinear form $Q[\mathbf{x}] = L(\mathbf{x}, \mathbf{x}) = (A\mathbf{x}, \mathbf{x})$.
 - Alternatively: A homogeneous polynomial of degree 2, i.e., a polynomial in x_1, \ldots, x_n with only ax_k^2 and cx_jx_k terms.
- There are infinitely many ways to write a quadratic form as $(A\mathbf{x}, \mathbf{x})$.
 - However, there is a unique representation $(A\mathbf{x}, \mathbf{x})$ where A is a (real) symmetric matrix.
- Quadratic form (on \mathbb{C}^n): A function of the form $Q[\mathbf{x}] = (A\mathbf{x}, \mathbf{x})$ where A is self-adjoint.
- Lemma 7.1.1: Let $(A\mathbf{x}, \mathbf{x})$ be real for all $\mathbf{x} \in \mathbb{C}^n$. Then $A = A^*$.
- To classify quadratic forms, consider the set of points $\mathbf{x} \in \mathbb{R}^n$ defined by $Q[\mathbf{x}] = 1$ for some quadratic form Q.
 - If the matrix of Q is diagonal, i.e., $Q[\mathbf{x}] = a_1 x_1^2 + \cdots + a_n x_n^2$, then the set of points can easily be visualized.
- The standard method of diagonalizing a quadratic form is change of variables.
- Orthogonal diagonalization.

- Let $Q[\mathbf{x}] = (A\mathbf{x}, \mathbf{x})$ in \mathbb{F}^n .
- Suppose $\mathbf{y} = S^{-1}\mathbf{x}$ where S is an invertible $n \times n$ matrix. Then

$$Q[\mathbf{x}] = Q[S\mathbf{y}] = (AS\mathbf{y}, S\mathbf{y}) = (S^*AS\mathbf{y}, \mathbf{y})$$

so in the new variables \mathbf{y} , the quadratic form has matrix S^*AS .

- Thus, we can let $A = UDU^*$, choose $D = U^*AU$ as our new (diagonal) matrix, and let this matrix act on the variables $\mathbf{y} = U^*\mathbf{x}$.
- Non-orthogonal diagonalization:
 - Completing the square:
 - Eliminate all $x_i x_j$ terms by completing the square. Then substitute in a y_k for each squared term.
 - Row/column operations:
 - Augment (A|I). Row reduce A to D. Then $I \to S^*$.

10/28:

- Silvester's Law of Inertia: For a Hermitian matrix A (i.e., for a quadratic form $Q[\mathbf{x}] = (A\mathbf{x}, \mathbf{x})$) and any of its diagonalizations $D = S^*AS$, the number of positive, negative, and zero diagonal entries of D depends only on A, but not on a particular choice of diagonalization.
- **Positive** (subspace $E \subset \mathbb{F}^n$ corresponding to A): A subspace E such that $(A\mathbf{x}, \mathbf{x}) > 0$ for all nonzero $\mathbf{x} \in E$. Also known as A-positive.
- Negative (subspace $E \subset \mathbb{F}^n$ corresponding to A): A subspace E such that $(A\mathbf{x}, \mathbf{x}) < 0$ for all nonzero $\mathbf{x} \in E$. Also known as A-negative.
- Neutral (subspace $E \subset \mathbb{F}^n$ corresponding to A): A subspace E such that $(A\mathbf{x}, \mathbf{x}) = 0$ for all nonzero $\mathbf{x} \in E$. Also known as A-neutral.
- Theorem 7.3.1: Let A be an $n \times n$ Hermitian matrix, and let $D = S^*AS$ be its diagonalization by an invertible matrix S. Then the number of positive (resp. negative) diagonal entries of D coincides with the maximal dimension of an A-positive (resp. A-negative) subspace.
- Lemma 7.3.2: Let $D = \text{diag}\{\lambda_1, \dots, \lambda_n\}$. Then the number of positive (resp. negative) diagonal entries of D coincides with the maximal dimension of a D-positive (resp. D-negative) subspace.
- Positive definite (quadratic form Q): A quadratic form Q such that $Q[\mathbf{x}] > 0$ for all $\mathbf{x} \neq \mathbf{0}$.
- Positive semidefinite (quadratic form Q): A quadratic form Q such that $Q[\mathbf{x}] \geq 0$ for all \mathbf{x} .
- Negative definite (quadratic form Q): A quadratic form Q such that $Q[\mathbf{x}] < 0$ for all $\mathbf{x} \neq \mathbf{0}$.
- Negative semidefinite (quadratic form Q): A quadratic form Q such that $Q[\mathbf{x}] \leq 0$ for all \mathbf{x} .
- Indefinite (quadratic form Q): A quadratic form Q for which there exist $\mathbf{x}_1, \mathbf{x}_2$ such that $Q[\mathbf{x}_1] > 0$ and $Q[\mathbf{x}_2] < 0$.
- Positive definite (Hermitian matrix A): A matrix A for which the corresponding quadratic form $Q[\mathbf{x}] = (A\mathbf{x}, \mathbf{x})$ is positive definite.
 - Positive semidefinite, negative definite, negative semidefinite, and indefinite Hermitian matrices are defined similarly.
- Theorem 7.4.1: Let $A = A^*$. Then
 - 1. A is positive definite iff all eigenvalues of A are positive.
 - 2. A is positive semidefinite iff all eigenvalues of A are non-negative.

- 3. A is negative definite iff all eigenvalues of A are negative.
- 4. A is negative semidefinite iff all eigenvalues of A are non-positive.
- 5. A is indefinite iff it has both positive and negative eigenvalues.
- Upper left submatrix (of A): A $k \times k$ matrix A_k composed of all entries of A from row (column) 1 through k in the same arrangement.
- Theorem 7.4.2 (Silvester's Criterion of Positivity): A matrix $A = A^*$ is positive definite if and only if $\det A_k > 0$ for all k = 1, ..., n.
 - To check if a matrix A is negative definite, check that the matrix -A is positive definite.
- Theorem 7.4.3 (Minimax characterization of eigenvalues): Let $A = A^*$ be an $n \times n$ matrix and let $\lambda_1 \ge \cdots \ge \lambda_n$ be its eigenvalues taken in decreasing order. Then

$$\lambda_k = \max_{E: \dim E = k} \min_{\mathbf{x} \in E: \|\mathbf{x}\| = 1} (A\mathbf{x}, \mathbf{x}) = \min_{F: \operatorname{codim} F = k-1} \max_{\mathbf{x} \in F: \|\mathbf{x}\| = 1} (A\mathbf{x}, \mathbf{x})$$

• Corollary 7.4.4 (Intertwining of eigenvalues): Let $A = A^* = \{a_{j,k}\}_{j,k=1}^n$ be a self-adjoint matrix and let $\tilde{A} = \{a_{j,k}\}_{j,k=1}^{n-1}$ be its submatrix of size $(n-1) \times (n-1)$. Let $\lambda_1, \ldots, \lambda_n$ and μ_1, \ldots, μ_{n-1} be the eigenvalues of A and \tilde{A} respectively, taken in decreasing order. Then

$$\lambda_1 \ge \mu_1 \ge \lambda_2 \ge \mu_2 \ge \cdots \ge \lambda_{n-1} \ge \mu_{n-1} \ge \lambda_n$$

Dual Spaces and Tensors

8.1 Notes

10/22: • Functional: A linear bounded map $L: H \to F$, where H is finite dimensional (equivalent to \mathbb{R}^n).

• Dual space: The set of bounded linear functionals on H. Denoted by H', H^* .

• If $l \leq p < \infty$, then

$$l^p = \left\{ (a_n)_{n \in \mathbb{N}} : \sum_{n=1}^{\infty} |a_n|^p < \infty \right\}$$

• Back to finite dimensions, $H' \approx \mathbb{R}^n$.

• Let $\mathbf{a}_1, \ldots, \mathbf{a}_n$ be a basis of H. Then $L\mathbf{x} = (L\mathbf{a}_1, \ldots, L\mathbf{a}_n) \approx \mathbb{R}^n$.

• Let $L((a_n)_{n\in\mathbb{N}}) = \sum_{n=1}^{\infty} a_n b_n$. Then $L((a_n)_{n\in\mathbb{N}})$ will be bounded if and only if $(b_n)_{n\in\mathbb{N}} \in l^q$ where $1 where <math>\frac{1}{q} + \frac{1}{p} = 1$.

• Young's inequality: The statement

$$ab \leq \frac{a^p}{p} + \frac{b^q}{q}$$

• We have $|\sum a_n b_n| \le ||a_n||_p ||b_n||_p$.

• Conclusion:

$$\sum \frac{|a_n||b_n|}{\|a_n\|_p \|b_n\|_q} = 1$$

• We can define H'', too. This contains linear functionals on H'.

• We know that $L(x) = \langle x, L \rangle = x(L)$. $x \in H''$.

• Riesz representation theorem: Let H have an inner product. $L \in H'$ if and only if there exists a unique $y \in H$ such that L(x) = (x, y).

- Gives us a way to identify all bounded linear functionals on H.

– In finite dimensions, L(x), where $x = \sum_{i=1}^{n} \alpha_i a_i$ gives us $L(x) = \sum_{i=1}^{n} \alpha_i L(a_i)$.

8.2 Chapter 8: Dual Spaces and Tensors

10/28: • Linear functionals are denoted by L.

- L is given by a $1 \times n$ matrix denoted by [L].
- The collection of all [L] (the dual space) is isomorphic to \mathbb{R}^n via $[L] \mapsto [L]^T$.
 - However, the objects are different: Let $[I]_{\mathcal{BA}}$ be the change of coordinates matrix in \mathbb{R}^n . We thus have that

$$[\mathbf{v}]_{\mathcal{B}} = [I]_{\mathcal{B}\mathcal{A}}[\mathbf{v}]_{\mathcal{A}}$$

but we also have that

$$[L]_{\mathcal{B}} = [L]_{\mathcal{A}}[I]_{\mathcal{A}\mathcal{B}}$$

so that

$$[L]_{\mathcal{B}}^T = ([L]_{\mathcal{A}}[I]_{\mathcal{A}\mathcal{B}})^T = [I]_{\mathcal{A}\mathcal{B}}^T[L]_{\mathcal{A}}^T$$

- Essentially, "if S is the change of coordinate matrix in X... then the change of coordinate matrix in the dual space X' is $(S^{-1})^T$ " (Treil, 2017, p. 219).
- Lemma 8.1.3: Let $\mathbf{v} \in V$. If $L(\mathbf{v}) = 0$ for all $L \in V'$, then $\mathbf{v} = \mathbf{0}$. As a corollary, if $L(\mathbf{v}_1) = L(\mathbf{v}_2)$ for all $L \in V'$, then $\mathbf{v}_1 = \mathbf{v}_2$.
- The second dual V'' is canonically (i.e., in a natural way) isomorphic to V.
- **Dual basis** (to $\mathbf{b}_1, \dots, \mathbf{b}_n \in V$): The system of vectors $\mathbf{b}'_1, \dots, \mathbf{b}'_n \in V'$ uniquely defined by the following equation. Also known as **biorthogonal basis**.

$$\mathbf{b}'_k(\mathbf{b}_i) = \delta_{ki}$$

- The k^{th} coordinate of a vector \mathbf{v} in a basis $\mathbf{b}_1, \dots, \mathbf{b}_n$ is $\mathbf{b}'_k(\mathbf{v})$.
 - This is a baby version of the abstract non-orthogonal Fourier decomposition of v.
- Theorem 8.2.1 (Riesz representation theorem): Let H be an inner product space. Given a linear functional L on H, there exists a unique vector $\mathbf{y} \in H$ such that

$$L(\mathbf{v}) = (\mathbf{v}, \mathbf{y})$$

for all $\mathbf{v} \in H$.

- If V is a real inner product space, we can define an isomorphism from V to V' by $\mathbf{y} \mapsto L_{\mathbf{v}} = (\mathbf{v}, \mathbf{y})$.
 - If V is complex, this function is not linear since if α is complex,

$$L_{\alpha \mathbf{v}}(\mathbf{v}) = (\mathbf{v}, \alpha \mathbf{y}) = \bar{\alpha}(\mathbf{v}, \mathbf{y}) = \bar{\alpha}L_{\mathbf{v}}(\mathbf{v})$$

- It follows by such a mapping that $\mathbf{b}'_k = \mathbf{b}_k$ for each k.
- Conjugate linear (transformation): A transformation T such that

$$T(\alpha \mathbf{x} + \beta \mathbf{v}) = \bar{\alpha} T \mathbf{x} + \bar{\beta} T \mathbf{v}$$

- It is customary to write outputs of linear functionals $L(\mathbf{v})$ in the form $\langle \mathbf{v}, L \rangle$.
 - This expression is linear in both arguments, unlike the inner product.
- Defines the dual transformation as the unique transformation such that

$$\langle A\mathbf{x}, \mathbf{y}' \rangle = \langle \mathbf{x}, A'\mathbf{y} \rangle$$

for all $\mathbf{x} \in X$, $\mathbf{y}' \in Y'$.

- It's matrix in the standard bases equals A^T .
- Annihilators are denoted by E^{\perp} here.
- \bullet Proposition 8.3.6: The annihilator of the annihilator of E equals E.
- Let $A:X\to Y$ be an operator acting from one vector space to another. Then
 - 1. $\ker A' = (\operatorname{range} A)^{\perp}$.
 - 2. $\ker A = (\operatorname{range} A')^{\perp}$.
 - 3. range $A = (\ker A')^{\perp}$.
 - 4. range $A' = (\ker A)^{\perp}$.

Advanced Spectral Theory

9.1 Notes

10/22:

- Let $p(z) = \sum_{i=0}^{n} a_i z^i$ be a polynomial. Let A be an $n \times n$ matrix. We let $p(A) = \sum_{i=0}^{n} a_i A^i$.
- Theorem: If A is an $n \times n$ and $p(\lambda) = \det(A \lambda I)$, then p(A) = 0.
 - We know that $p(\lambda) = a(z \lambda_1) \cdots (z \lambda_n)$ where $\lambda_1, \dots, \lambda_n$ are the eigenvalues.
 - Thus $p(A) = a(A \lambda_1 I) \cdots (A \lambda_n I)$.
 - If you are in \mathbb{R}^n and have this property, you can factorize your matrix.
 - Thus, $p(A)\mathbf{x} = \mathbf{0}$ since \mathbf{x} can be decomposed into a linear combination of eigenvectors of A, which will be taken to 0 one by one by the terms of p(A).
- $\sigma(B) = \{\text{eigenvalues of } B\}$ is known as the **spectrum** of B.
- If p is an arbitrary polynomial and A is $n \times n$, then μ is an eigenvalue of p(A) if and only if $\mu = p(\lambda)$ where λ is an eigenvalue of A. In essence, $\sigma(p(A)) = p(\sigma(A))$.
- Chapter 9 will not be on the exam. We don't have to know the generalization to infinite dimensional spaces.

10/25:

- If A is an $n \times n$ square matrix and $p(\lambda) = \det(A \lambda I)$, then p(A) = 0.
 - Proof: WLOG, let A be an upper triangular matrix with diagonal entries equal to the eigenvalues.
 - Think of $p(z) = (-1)^n (z \lambda_1) \cdots (z \lambda_n)$.
 - Thus, $p(A) = (-1)^n (A \lambda_1 I) \cdots (A \lambda_n I)$.
 - WTS: $p(A)\mathbf{x} = 0$ for all $\mathbf{x} \in V$.
 - Let $E_k = \operatorname{span}(e_1, \dots, e_k)$ be the span of the first k eigenvectors of A, where e_1, \dots, e_n is a standard basis in \mathbb{C}^n .
 - A triangular implies $AE_k \subset E_k$. Thus, $(A \lambda I)E_k \subset E_k$, so E_k is invariant under $A \lambda I$ for all λ .
 - If we apply $A \lambda_k I$ to a vector in E_k , we are left with a vector in E_{k-1} .
 - Thus, if we apply $\prod_{k=1}^{n} (A \lambda_k I) = p(A)$ to any vector in $E_n = V$, we will kill it piece by piece down to zero.
- Let A be a square $n \times n$ matrix. Then p an arbitrary polynomial implies $\sigma(p(A)) = p(\sigma(A))$. (Any eigenvalue μ of p(A) is $\mu = p(\lambda)$, where λ is an eigenvalue of A.)
 - Shows that polynomials of operators commute.

- Proof: Let λ be an eigenvalue of A. We want to show that $p(\lambda)$ is an eigenvalue of p(A). This is obvious since $A\mathbf{x} = \lambda \mathbf{x}$ for some \mathbf{x} , so $A^k \mathbf{x} = \lambda^k \mathbf{x}$, so in particular, $p(A)\mathbf{x} = p(\lambda)\mathbf{x}$.
- On the other hand, if μ is an eigenvalue of p(A), we want to show that there exists $\lambda \in \sigma(A)$ such that $\mu = p(\lambda)$.
- Consider $q(z) = p(z) \mu$. Then $q(A) = p(A) \mu I$. Since μ is an eigenvalue of p(A), q(A) is not invertible.
- Thus, $q(z) = (-1)^n (z z_1) \cdots (z z_n)$ and $q(A) = (-1)^k (A z_1 I) \cdots (A z_k I)$.
- But q(A) is not invertible, so one of the $A z_k I$ is not invertible. Take z_k such that $A z_k I$ is not invertible. Then $z_k \in \sigma(A)$. It follows that $q(z_k) = p(z_k) \mu = \sigma$.
- If A is $n \times n$, $\lambda_1, \ldots, \lambda_n$ are its eigenvalues, p is a polynomial, then p(A) is invertible if and only if $p(\lambda_k) \neq 0$ for each $k = 1, \ldots, n$.
 - This is an immediate corollary to the previous result.
- We now build up to the **generalized eigenspace**, which is related to some "geometric" properties of the algebraic multiplicity of an eigenvalue.
- If $A: V \to V$ is a linear operator and $E \subset V$ is a subspace, E is A-invariant if $AE \subset E$.
- Facts:
 - If E is A-invariant, E is A^k -invariant.
 - Thus, E is p(A)-invariant.
- Consider the restriction map $A|_E$.
- A has a block-diagonalized matrix where each block corresponds to the generalized eigenvectors of a generalized eigenvalue of A.
 - Let E_1, \ldots, E_r be a basis of invariant subspaces.
 - Let $A_k = A|_{E_k}$. Then the A_k 's act independently of each other.
- Generalized eigenvector (of A): A vector \mathbf{v} corresponding to an eigenvalue λ if there exists $k \geq 1$ such that $(A \lambda I)^k \mathbf{v} = \mathbf{0}$.
- Generalized eigenspace: The set E_{λ} of all of the generalized eigenvectors of λ . Given by

$$E_k = \bigcup_{k \ge 1} \ker(A - \lambda I)^k$$

- $-E_{\lambda}$ is a linear subspace of V.
- **Degree** (of λ): The smallest number k such that increasing k any more does not add further vectors to the generalized eigenspace. Denoted by $d(\lambda)$. Also known as **depth**.
 - Symbolically, $d(\lambda)$ is the smallest number such that

$$E_{\lambda} = \bigcup_{k=1}^{d(\lambda)} \ker(A - \lambda I)^k$$

- Start working through the first 25 problems of Rudin (1976) (his metric spaces problems).
- 10/27: Jordan form.
 - \bullet Reviews build up to generalized eigenvectors.

- Theorem: If $\sigma(A) = \{\lambda_1, \dots, \lambda_n\}$ and E_1, \dots, E_n are the corresponding generalized eigenspaces, then E_1, \dots, E_n is a basis of subspaces of U, i.e., $V = \bigoplus_k E_k$.
- Corollary: $A: V \to V$ can be represented as A = D + N where D is diagonalizable and N is nilpotent and ND = DN.
 - Proof: Consider the basis of generalized eigenspaces known to exist from the theorem. Then $A = \text{diag}\{A_1, \dots, A_r\}$.
 - Let

$$N_k = A_k - \lambda_k I_{E_k}$$

This is nilpotent.

- Then let

$$D = \operatorname{diag}\{\lambda_1 I_{E_1}, \dots, \lambda_n I_{E_n}\}$$

- These two matrices satisfy the necessary properties.
- Let $\dot{\mathbf{x}} = A\mathbf{x}$.
 - Let $\mathbf{x}(t) = e^{tA}$, where

$$e^{tA} = \sum \frac{(tA)^k}{k!}$$

- $-\|e^{tA}\| \le \sum \frac{\|A^k\|}{k!} = \sum \frac{\|A\|^k}{k!}.$
- Let p be a polynomial of degree k. Then

$$p(a+x) = \sum_{k=0}^{d} \frac{p^{(k)}(a)}{k!} x^k$$

- If A = D + N, then...
- Nilpotent operators:
 - Let $A = \operatorname{diag}\{A_1, \dots, A_r\}$.
 - We know that $A_k = \lambda_k I_{E_k} + N_k$ for each k.
 - Every nilpotent N can be written in the form

$$\begin{pmatrix} 0 & 1 & & 0 \\ & \ddots & \ddots & \\ & & \ddots & 1 \\ 0 & & & 0 \end{pmatrix}$$

- The exam is long but not that hard. The only question is will you do good or very good.
 - Revise the previous two homeworks, especially the last two.
 - No justification for any of the true/false questions. Just circle T or F.
 - There are four problems. One is T/F (with multiple subparts); the other three are problem problems (with subparts, too).
 - Some of the questions will take you 2 seconds. Some you've already seen in the PSets.
 - The exam is supposed to be boring.
 - Calculators?
 - No calculators needed. Calculators are for chem/physics exams.

- Not a lot of computation.
- 50 minutes.
- Chloe will be proctoring.
- Remember the determinant of "special" matrices.
 - $|\det U| = 1$ if U is unitary.
 - $-\det A = \pm 1$ if A is orthogonal.
 - Make a list of matrix types that are automatically diagonalizable.
 - Determinant is the product of the eigenvalues.
 - Determinant of A is equal to the conjugate of the determinant of A^* .
- Most of the exercises use the inner product.
 - Whenever you had something to prove about eigenvalues or eigenbasis, you went through diagonalization or SVD or the inner product or polar decomposition.
 - Proving eigenvalues of self-adjoint matrices are real w/ the inner product.
- Eigenvalues/eigenvectors of a projection.
 - It's implied that it's asking you the multiplicities!!!
- Know useful facts but have an idea how to prove them as well.
- Recommends against shorthanding in the exams.
- Not grading on clarity (since the exam is long).
- Max and min are for when you're sure something will be attained. Otherwise use sup and inf.

9.2 Chapter 9: Advanced Spectral Theory

- Theorem 9.1.1 (Cayley-Hamilton): If p is the characteristic polynomial of A, p(A) = 0.
 - Theorem 9.2.1 (Spectral Mapping Theorem): For a square matrix A and an arbitrary polynomial p, $\sigma(p(A)) = p(\sigma(A))$. In other words, μ is an eigenvalue of p(A) if and only if $\mu = p(\lambda)$ for some eigenvalue λ of A.
 - Corollary 9.2.2: Let A be a square matrix with eigenvalues $\lambda_1, \ldots, \lambda_n$ and let p be a polynomial. Then p(A) is invertible iff $p(\lambda_k) \neq 0$ for all $k = 1, \ldots, n$.
 - Algebraic multiplicity is the dimension of the corresponding generalized eigenspace.

Labalme 42

Part II Point Set Topology of Metric Spaces

The Real and Complex Number Systems

1.1 Notes

- 11/1: Spent a lot of time trying to cheer us up regarding the midterm.
 - There may be some true/false on linear algebra on the final.
 - Facts:
 - 1. $\sqrt{2}$ is irrational.
 - 2. Archimedes principle: If x > 0 and $y \in \mathbb{R}$, then there exists n such that nx > y.
 - 3. If x > y, then there exists $q \in \mathbb{Q}$ such that x > q > y.

1.2 Chapter 1: The Real and Complex Number Systems

From Rudin (1976).

11/6:

- Rudin (1976) presents several interesting proofs throughout this section that may be of interest later by means of their divergence from the ones with which I am familiar.
- Least-upper-bound property: The property pertaining to a set S that if $E \subset S$, $E \neq \emptyset$, and E is bounded above, then $\sup E \in S$.
 - For example, $\mathbb Q$ does not have the least-upper-bound property.
 - The **greatest-lower-bound property** is analogously defined.
- Theorem: Suppose S is an ordered set with the least-upper-bound property, $B \subset S$ is nonempty, and B is bounded below. Let L be the set of all lower bounds of B. Then $\alpha = \sup L$ exists in S, and $\alpha = \inf B$. In particular, $\inf B$ exists in S.
 - Essentially, this theorem states that any set that satisfies the least-upper-bound property satisfies
 the greatest lower bound property.
- Existence theorem: There exists an ordered field \mathbb{R} which has the least-upper-bound property. Moreover, \mathbb{R} contains \mathbb{Q} as a subfield.
 - The second statement implies that the operations of addition and multiplication on \mathbb{R} , when applied to \mathbb{Q} , coincide with the operations of addition and multiplication on \mathbb{Q} .

- Archimedean property (of \mathbb{R}): If $x \in \mathbb{R}$, $y \in \mathbb{R}$, and x > 0, then there is a positive integer n such that nx > y.
- Rudin (1976) proves several theorems about the real numbers from the least-upper-bound property as opposed to the traditional construction of the real numbers.
- Introduces the decimal system.
- Finite real number system: That which has been defined thus far.
- Extended real number system: The set $\mathbb{R} \cup \{+\infty, -\infty\}$ where $+\infty, -\infty$ obey the expected properties (supremum [resp. infimum] of every set, $x + \infty = \infty$, etc.).
- Defines the complex field axiomatically with complex numbers in the form (a,b) for $a,b \in \mathbb{R}$.
 - Notes that the real numbers form a subfield of the complex field.
 - Defines i = (0,1), proves $i^2 = -1$, proves a + bi = (a,b).
- Schwarz inequality: If a_1, \ldots, a_n and b_1, \ldots, b_n are complex numbers, then

$$\left| \sum_{j=1}^{n} a_j \bar{b}_j \right|^2 \le \sum_{j=1}^{n} |a_j|^2 \sum_{j=1}^{n} |b_j|^2$$

• Euclidean k-space: The vector space \mathbb{R}^k over the real field.

Basic Topology

2.1 Notes

11/1: • Equivalence relationships are denoted $A \sim B$.

- These are...
 - Reflexive $(A \sim A)$.
 - Symmetric $(A \sim B \iff B \sim A)$.
 - Transitive $(A \sim B \& B \sim C \Longrightarrow A \sim C)$.
- Equivalence relations give rise to equivalence classes.
- Countable (set A): A set A such that $A \sim \mathbb{N}$, in the sense that there exists a one-to-one and onto map from $\mathbb{N} \to A$.
 - Alternatively, A can be written in the form $A = \{f(n) : n \in \mathbb{N}\}.$
- Finite countable vs. infinite countable (see Rudin (1976)).
- N denotes the natural numbers.
- \mathbb{N}_0 denotes the natural numbers including 0.
- \mathbb{Z} denotes the integers.
- We know that $\mathbb{N} \sim \mathbb{Z}$: Let $f: \mathbb{N} \to \mathbb{Z}$ be defined by

$$f(n) = \begin{cases} \frac{n}{2} & n \text{ even} \\ \frac{n-1}{2} & n \text{ odd} \end{cases}$$

- More facts.
 - 1. Every subset of a countable set is countable.
 - 2. Unions of countable sets are countable.
 - If the sets E_n for some finite list of numbers are countable, then $\bigcup_n E_n$ is countable.
 - Soug goes over the diagonalization method of counting.
 - 3. n-fold Cartesian products of countable sets are countable (we induct on n).
 - If A is countable and B is countable, then $A \times B$ is countable.
 - If A is finite and to each $\alpha \in A$ we assign a countable set E_{α} , $\otimes_{\alpha \in A} E_{\alpha}$ is countable.
- Metric space: A space X along with a metric $d: X \times X \to [0, \infty)$ such that

- $-d(x,y) > 0 \text{ iff } x \neq y, \text{ and } d(x,x) = 0 \text{ iff } x = 0.$
- d(x,y) = d(y,x).
- $d(x,y) \le d(x,z) + d(z,y).$
- Example (\mathbb{R}^n) :
 - We may define d by

$$d(x,y) = \sqrt{\sum (x_i - y_i)^2}$$

- We can also define the p-metrics (recall normed spaces) with p where 2 is.
- Example $(X_p = \{f : Y \to \mathbb{R} : 1 \le p < \infty, \int_Y |f|^p dy < \infty\})$:
 - This is ℓ_p .
 - Define

$$||f - g||_p = \left[\int_Y |f - g|^p \, \mathrm{d}y \right]^{1/p}$$

- Convergence: $x_n \to x \iff d(x_n, x) \to 0$.
- Neighborhood: The set of all points a distance less than r away from p. Denoted by $N_r(p)$. Given by

$$N_r(p) = \{ q \in X : d(p,q) < r \}$$

- **Limit point** (of *E*): A point *p* such that every neighborhood of *p* intersects *E* at a point other than *p*. Also known as **accumulation point**.
 - Symbolically,

$$N_r(p) \cap (E \setminus \{p\}) \neq \emptyset$$

for all r > 0.

- Isolated point (of E): A point p such that $p \in E$ and p is not a limit point of E.
- Closed (set E): A set E that contains all of its limit points.
- Interior (point p): A point p such that there exists $N_r(p) \subset E$.
- Open (set E): A set E, all points of which are interior points.
- **Perfect** (set E): A set E that is closed and every point of E is a limit point of E.
- Bounded (set E): There exists a number M and a $y \in X$ such that $E \subset \{p : d(p,y) \le M\}$.
- Dense (set E in X): A set E such that every point of X is a limit point of E or a point of E, itself.
- 11/3: Every neighborhood is an open set.
 - If p is a limit point of E, every neighborhood of p contains infinitely many points of E.
 - Thus, a finite set cannot have a limit point.
 - Prove by contradiction: Suppose there is a neighborhood that contains only finitely many points of E. Then the neighborhood with radius smaller than the distance to the closest point does not contain any points of E, a contradiction.
 - E is open iff $E^{c[1]}$ is closed.
 - Assume E^c closed. If $p \in E$, then p is not a limit point of E^c . It follows that there exists a neighborhood of p that is entirely contained within E, so p is interior, as desired.

 $^{^{1}}$ The complement of E.

- Suppose E is open. Let p be any limit point of E^c . Then $p \in E^c$.
- F is closed iff F^c is open.
- If $(G_{\alpha})_{\alpha \in A}$ is a family of open sets in X, then the union is open.
 - Let $p \in \bigcup_{\alpha \in A} G_{\alpha}$. Then $p \in G_{\alpha}$ for some $\alpha \in A$. It follows that p is an interior point of G_{α} , so thus an interior point of the union of G_{α} with everything else.
- Finite intersections of open sets are open.
 - In the infinite case $\bigcap_{n\in\mathbb{N}}(-1/n,1/n)=\{0\}$, an intersection of infinitely many open sets is closed.
 - However, in the finite case, just consider the neighborhood with the smallest radius and take this
 one.
- The intersection of closed sets is closed.
- The union of finitely many closed sets is closed.
 - These follow from the previous two by De Morgan's rule.
- Let $\bar{E} = E \cup E'$ where E' is the set of limit points of E.
- Let X be a metric space and $E \subset X$. Then
 - 1. \bar{E} is closed.
 - WTS: \bar{E}^c is open. Let $p \in \bar{E}^c$. Then p is neither in E nor is it a limit point of E. Thus, there exists a neighborhood of \bar{E}^c containing entirely points of \bar{E}^c . Therefore, \bar{E}^c is open, so \bar{E} is closed.
 - 2. $E = \bar{E}$ iff E is closed.
 - Think $p \in \bigcap G_{\alpha}$?
 - 3. $\bar{E} \subset F$ for any closed $F \supset E$.
 - If $E \subset F$, then any limit point of E will be a limit point of F. Thus, $E' \subset F'$. Then $\bar{E} = E \cup E' \subset F \cup F' = \bar{F} = F$ where the last equality holds because F is closed.
- Types of sets.

	Closed	Open	Perfect	Bounded
	N	Y	N	Y
$\{z\in\mathbb{Q}: z \leq 1\}$	Y	N	Y	Y
Nonempty finite set	Y	N	N	Y
\mathbb{Z}	Y	N	N	N
$\{1/n:n\in\mathbb{N}\}$	N	N	N	Y
\mathbb{R}^2	Y	Y	Y	N
(a,b)	N	?	N	Y

Table 2.1: Types of sets.

- Relatively open (set E to Y): A set $E \subset Y \subset X$ such that if $p \in E$, then there exists a Y-neighborhood of E contained in E.
- Let $N_r^X(p) = \{y \in X : d(y,p) < r\}$ be a neighborhood of p in X, and let $N_r^Y(p) = \{y \in Y : d(y,p) < r\}$ be a neighborhood of p in Y. Then $N_r^Y(p) = N_r^X(p) \cap Y$.

- E is open relative to Y iff $E = G \cap Y$ where G is open relative to X.
- Introduces the supremum.
- If $E \subset \mathbb{R}$, $E \neq \emptyset$, and E is bounded above, sup $E < \infty$.
- Let $y = \sup E$. Then $y \in \bar{E}$.
- There exists a sequence $a_n \in A$ such that $a_n \to x = \sup A$.
- A is compact iff any open cover of the set has a finite subcover.
- Study and *know* all of these proofs.
- 11/5: Compactness: Defines compactness in terms of open covers.
 - Finite sets are compact.
 - Compactness is "absolute" (i.e., it is not a relative property like openness).
 - If $K \subset Y \subset X$, then K is compact relative to X iff K is compact relative to Y.
 - -V is open relative to Y iff $V=G\cap Y$ where G is open relative to X.
 - Compact implies closed.
 - We will show K compact implies K^c open.
 - WTS: For all $p \in K^c$, there exists $N_r(p) \subset K^c$ sub that $N_r(p) \cap K = \emptyset$.
 - A closed subset of a compact set is compact.
 - Let K be compact and let $F \subset K$ be closed.
 - Take any open cover of F. Extend it to an open cover of K. Take the finite subcover of K. Naturally, this finite subcover is also a finite cover of $F \subset K$.
 - F closed, K compact implies $F \cap K$ compact.
 - If $(K_{\alpha})_{\alpha \in A}$ is compact in X with finite intersection property (every intersection of any finite number of these sets is nonempty), then $\bigcap_{\alpha \in A} K_{\alpha} \neq \emptyset$.
 - Argue by contradiction.
 - Let $G_{\alpha} = K_{\alpha}^{c}$.
 - Assume the intersection is empty. Assume WLOG that no point of K_1 is in any of the other K_{α} 's.
 - Then $\{G_{\alpha}\}_{{\alpha}\in A}$ be an open cover of K_1 .
 - K_1 compact implies there is a finite subcover $G_{\alpha_1}, \ldots, G_{\alpha_n}$. Then $K \subset G_{\alpha_1} \cup \cdots \cup G_{\alpha_n}$. This implies that $K_1 \cap K_{\alpha_1} \cap \cdots \cap K_{\alpha_n} = \emptyset$, a contradiction.
 - Let E be an infinite subset of a compact K. Then E has a limit point in K.
 - Argue by contradiction.
 - Suppose for all $p \in K$, there exists $N_r(p)$ such that $N_r(p) \cap E = \{p\}$.
 - Consider the set $\{N_r(p): p \in K\}$. This is an open cover of K. Thus, there exists a finite subcover of it. But since $E \subset K \subset N_{r_1}(p_1) \cup \cdots \cup N_{r_n}(p_n) = \{p_1\} \cup \cdots \cup \{p_n\}$, E is finite, a contradiction.
 - 2-cell (in \mathbb{R}^2): A set that is the Cartesian product of two closed intervals.
 - Generalizes to k-cells.
 - Let $I_n = [a_n, b_n] \subset \mathbb{R}$ such that $I_{n+1} \subset I_n$. Then $\bigcap_{n=1}^{\infty} I_n \neq \emptyset$.

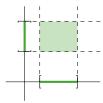


Figure 2.1: 2-cell.

- Let I_k be a k-cell in \mathbb{R}^k such that $I_k \supset I_{k+1}$. Then $\bigcap_k I_k \neq \emptyset$.
- We know that $a_m \leq a_{m+n} \leq b_{m+n} \leq b_m$, so $\sup a_n \in \bigcap I_n$.
- Every k-cell is compact.

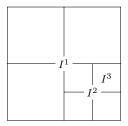


Figure 2.2: k-cells are compact.

- Argue by contradiction.
- Consider an open cover of the k-cell I^1 . If it has a finite subcover, we're done. So suppose we have an open cover that doesn't have a finite subcover. Split the k-cell into 2^k chunks. At least one of the chunks I^2 must not have a finite subcover.
- Split that one into 2^k chunks. At least one of the chunks I^3 must not have a finite subcover.
- Continue.
- Thus, we have a decreasing family of k-cells, so by the previous result, their $\bigcap I^n \neq \emptyset$.
- Let $x \in \bigcap I^n$. Then the...
- Heine-Borel theorem: Let $E \subset \mathbb{R}^k$. Then TFAE^[2]
 - 1. E is closed and bounded.
 - 2. E is compact.
 - 3. Every infinite subset of E has a limit point in E.
 - (1 \Rightarrow 2) E closed and bounded implies E is a closed subset of some I_k , so it's compact.
 - $-(2 \Rightarrow 3)$ Already done.
 - $-(3 \Rightarrow 1)$
 - \blacksquare Suppose E not bounded. Then there is an infinite sequence of points in E that never converges. Contradiction.
 - Suppose E is not closed. Then there exists a sequence of points in E which "converges" to an $x_0 \notin E$.

 $^{^2}$ The following are equivalent.

2.2 Chapter 2: Basic Topology

From Rudin (1976).

- Countable (set A): A set A that is in bijective correspondence with the set of all positive integers.

 Also known as enumerable, denumerable.
 - At most countable (set A): A set A that is finite or countable.
 - An alternative definition of an **infinite** set would be a set that is equivalent to one of its proper subsets.
 - ullet Theorem: Let A be the set of all sequences whose elements are the digits 0 and 1. This set A is uncountable.

Proof. Let $E = \{s_1, s_2, \dots\}$ be an arbitrary countable subset of A, where each s_j is a sequence whose elements are the digits 0 and 1. Let s be the sequence, the n^{th} term of which is the opposite of the n^{th} term of s_n (i.e., if the n^{th} term of s_n is 0, we set the n^{th} term of s equal to 1). This guarantees that s is distinct from each of the s_j , i.e., that $s \notin E$. It follows that $E \subsetneq A$, i.e., that every countable subset of A is a proper subset of A. Therefore, A must be uncountable (for otherwise A would be a proper subset of A, a contradiction).

- The idea of this proof is called **Cantor's diagonalization process**.
- Since every real number can be represented as a binary sequence of numbers, i.e., $A \sim \mathbb{R}$, the reals are uncountable.
- Metric space: A set X such that with any two points $p, q \in X$, there is associated a real number d(p,q) such that
 - 1. d(p,q) > 0 if $p \neq q$; d(p,p) = 0.
 - 2. d(p,q) = d(q,p).
 - 3. $d(p,q) \le d(p,r) + d(r,q)$ for any $r \in X$.
- Distance (from $p \in X$ to $q \in X$, X a metric space): The real number d(p,q).
- Distance function: A function $d: X \times X \to \mathbb{R}$ that sends $(p,q) \mapsto d(p,q)$. Also known as metric.
- Every subset of a metric space is a metric space in its own right under the same distance function.
- Segment (from a to b): The set of all real numbers x such that a < x < b. Denoted by (a, b).
- Interval (from a to b): The set of all real numbers x such that $a \le x \le b$. Denoted by [a, b].
- **k-cell**: The set of all points $\mathbf{x} = (x_1, \dots, x_k) \in \mathbb{R}^k$ whose coordinates satisfy the inequalities $a_i \leq x_i \leq b_i$ where $a_i < b_i$ for each $1 \leq i \leq k$.
 - $-\,$ Note that a 1-cell is an interval and a 2-cell is a rectangle.
- Convex (set E): A subset E of \mathbb{R}^k such that

$$\lambda \mathbf{x} + (1 - \lambda) \mathbf{y} \in E$$

for all $\mathbf{x}, \mathbf{y} \in E$ and $0 < \lambda < 1$.

- Balls and k-cells are both convex.
- The segment (a,b) is open as a subset of \mathbb{R}^1 , but not open as a subset of \mathbb{R}^2 .
- Since compactness is not relative, while it makes no sense to talk about *open* or *closed* metric spaces, it does make sense to talk about *compact* metric spaces.

- Weierstrass theorem: Every bounded infinite subset of \mathbb{R}^k has a limit point in \mathbb{R}^k .
- Theorem: Let P be a nonempty perfect set in \mathbb{R}^k . Then P is uncountable.

Proof. Since P is nonempty and perfect, there exists a limit point of P. It follows that P is infinite.

Now suppose for the sake of contradiction that P is countable, and denote the elements of P by $\mathbf{x}_1, \mathbf{x}_2, \ldots$. We now construct a sequence $\{V_n\}$ of neighborhoods, as follows. Let $V_1 = N_r(\mathbf{x}_1)$. Clearly, $V_1 \subset P$ since $\mathbf{x}_1 \in P$. It follows that since V_1 is a neighborhood that V_1 contains infinitely many points of P. Now suppose inductively that V_n has been constructed. Thus, by analogous conditions to those on V_1 , we may let V_{n+1} be a neighborhood such that (i) $\bar{V}_{n+1} \subset V_n$, (ii) $\mathbf{x}_n \notin \bar{V}_{n+1}$, and (iii) $V_{n+1} \cap P$ is nonempty. By (iii), we can continue on to construct V_{n+2} , and so on and so forth.

Let $K_n = \bar{V}_n \cap P$. Since \bar{V}_n is closed and bounded, \bar{V}_n is compact. Additionally, since $\mathbf{x}_n \notin K_{n+1}$ for each n, no point of P lies in $\bigcap_{1}^{\infty} K_n$. Thus, since each $K_n \subset P$, $\bigcap_{1}^{\infty} K_n$ is empty. But this contradicts our previous result that since each K_n is nonempty, compact, and such that $K_n \supset K_{n+1}$, $\bigcap_{1}^{\infty} K_n$ is nonempty.

- Corollary: Every interval [a, b] is uncountable. In particular, \mathbb{R} is uncountable.
- Cantor set: The set resulting from the following construction. Let $E_0 = [0, 1]$. Remove the segment (1/3, 2/3), so that $E_1 = [0, 1/3] \cup [1/3, 2/3]$. Now remove the middle third of these two intervals to create E_2 . Continue on indefinitely.
 - This is a perfect set in \mathbb{R}^1 which contains no segment.
- Separated (sets A, B): Two subsets A, B of a metric space X such that $A \cap \overline{B}$ and $\overline{A} \cap B$ are empty.
- Connected (set E): A set E that is not the union of two nonempty separated sets.
- Separated sets are disjoint, but disjoint sets are not necessarily separated (consider [0,1] and (1,2)).
- A subset E of \mathbb{R}^1 is connected if and only if it has the following property: If $x, y \in E$ and x < z < y, then $z \in E$.

Numerical Sequences and Series

3.1 Chapter 3: Numerical Sequences and Series

From Rudin (1976).

11/7: • Convergence of sequences is relative.

- For example, the sequence 1/n for $n=1,2,\ldots$ converges in \mathbb{R} , but not in $(0,\infty)$.
- Range (of $\{p_n\}$): The set of all points p_n .
 - This definition squares nicely with the formal definition of a sequence as a function p defined on \mathbb{N} .
- Theorem 3.6a: If $\{p_n\}$ is a sequence in a compact metric space X, then some subsequence of $\{p_n\}$ converges to a point of X.
- Theorem 3.7: The subsequential limits of a sequence $\{p_n\}$ in a metric space X form a closed subset of X.
- **Diameter** (of E): The supremum of the set

$$S = \{d(p,q): p,q \in E\}$$

where E is a nonempty subset of a metric space X. Denoted by diam E.

- Theorem 3.10:
 - (a) If \bar{E} is the closure of a set E in a metric space X, then

$$\dim \bar{E} = \dim E$$

- (b) If K_n is a sequence of compact sets in X such that $K_n \supset K_{n+1}$ (n = 1, 2, 3, ...) and if $\lim_{n\to\infty} \operatorname{diam} K_n = 0$, then $\bigcap_{1}^{\infty} K_n$ consists of exactly one point.
- Complete (metric space): A metric space in which every Cauchy sequence converges.
- All compact metric spaces and all Euclidean spaces are complete.
 - The metric space $(\mathbb{Q}, |x-y|)$ is not complete.
- Monotonically increasing (sequence $\{s_n\}$): A sequence $\{s_n\}$ of real numbers such that $s_n \leq s_{n+1}$ for each $n \in \mathbb{N}$.
- Monotonically decreasing (sequence $\{s_n\}$): A sequence $\{s_n\}$ of real numbers such that $s_n \geq s_{n+1}$ for each $n \in \mathbb{N}$.

- Monotonic sequences: The class of all sequences that are either monotonically increasing or monotonically decreasing.
- Upper limit (of $\{s_n\}$): The supremum of the set E of all subsequential limits of $\{s_n\}$. Denoted by s^* , $\limsup_{n\to\infty} s_n$.
- Lower limit (of $\{s_n\}$): The infimum of the set E of all subsequential limits of $\{s_n\}$. Denoted by s_* , $\lim \inf_{n\to\infty} s_n$.
- Theorem 3.17: Let $\{s_n\}$ be a sequence of real numbers. Then s^* has (and is the only number to have both of) the following two properties.
 - (a) $s^* \in E$.
 - (b) If $x > s^*$, then there is an integer N such that $n \ge N$ implies $s_n < x$.

An analogous result holds for s_* .

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