# Math 307 Study Guide

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# Chapter 1: Linear equations

### 1.1 Solving linear equations

- Vector norms
  - p-norm

$$||x||_p = (|x_1|^p + \dots + |x_n|^p)^{\frac{1}{p}}$$

- 1-norm

$$||x||_1 = |x_1| + \dots + |x_n|$$

- 2-norm

$$||x||_2 = \sqrt{|x_1|^2 + \dots + |x_n|^2}$$

 $-\infty$ -norm

$$||x||_{\infty} = \max\{|x_1|, ..., |x_n|\}$$

•  $\|\mathbf{x}\|$  is a norm iff the following properties hold

$$-\|\mathbf{x}\| \ge 0$$

$$-\|\mathbf{x}\| = 0 \text{ iff } \mathbf{x} = 0$$

$$-\|s\mathbf{v}\| = |s| \cdot \|\mathbf{v}\|$$

$$-\|\mathbf{u} + \mathbf{v}\| \le \|\mathbf{u}\| + \|\mathbf{v}\|$$

- Matrix norms
  - Hilbert-Schmidt (Frobenius) norm

$$||A||_{HS} = \left(\sum_{j=1}^{n} \sum_{i=1}^{m} |a_{ij}|^2\right)^{\frac{1}{2}}$$

- Matrix (operator) norm

$$||A|| = \max_{x \neq 0} \frac{||Ax||_2}{||x||_2} = \max_{||z||=1} ||Az||$$

If A is diagonal, then  $||A|| = \max\{|\lambda_1|, ..., |\lambda_n|\}$ 

• Condition number

$$k(A) = \left\|A^{-1}\right\| \|A\| = \frac{\text{largest stretch factor}}{\text{smallest stretch factor}}$$

If A is diagonal and invertible, then

$$k(A) = \frac{\max\{|\lambda_1|, ..., |\lambda_n|\}}{\min\{|\lambda_1|, ..., |\lambda_n|\}}$$

### 1.2 Interpolation

Given a set of data points, we want to fit a function exactly.

• Lagrange interpolation: Fit to  $p(x) = a_1 x^{n-1} + ... + a_{n-1} x + a_n$ .

$$\begin{bmatrix} x_1^{n-1} & \dots & x_1 & 1 \\ \vdots & & \vdots & \vdots \\ x_n^{n-1} & \dots & x_n & 1 \end{bmatrix} \begin{bmatrix} a_1 \\ \vdots \\ a_n \end{bmatrix} = \begin{bmatrix} y_1 \\ \vdots \\ y_n \end{bmatrix}$$

• Cubic spline interpolation: Fit to piecewise function

$$f(x) = \begin{cases} p_1(x) & x_1 \le x \le x_2 \\ \vdots & \\ p_{n-1}(x) & x_{n-1} \le x \le x_n \end{cases}$$
$$p_j(x) = A_j x^3 + B_j x^2 + C_j x + D_j$$

We impose conditions so that f(x) is continuous at all points and smooth.

### 1.3 Finite difference approximations

Given a differential equation f''(x) + q(x)f(x) = r(x) and initial conditions  $f(x_0) = A$ ,  $f(x_N) = B$ , we want to find an N + 1-point approximation of f(x).

$$\left(\begin{bmatrix}
1 & 0 & 0 & 0 & \dots & 0 & 0 \\
1 & -2 & 1 & 0 & \dots & 0 & 0 \\
0 & 1 & -2 & 1 & \dots & 0 & 0 \\
0 & 0 & 1 & -2 & \dots & 0 & 0 \\
\vdots & \vdots & \vdots & \vdots & \ddots & \vdots & \vdots \\
0 & 0 & 0 & 0 & \dots & -2 & 1 \\
0 & 0 & 0 & 0 & \dots & 0 & 1
\end{bmatrix} + (\Delta x)^2 \begin{bmatrix}
0 & 0 & 0 & \dots & 0 & 0 \\
0 & q_1 & 0 & \dots & 0 & 0 \\
0 & 0 & q_2 & \dots & 0 & 0 \\
\vdots & \vdots & \vdots & \ddots & \vdots & \vdots \\
0 & 0 & 0 & \dots & q_{N-1} & 0 \\
0 & 0 & 0 & \dots & 0 & 0
\end{bmatrix} \begin{bmatrix}
f_0 \\
f_1 \\
f_2 \\
\vdots \\
f_{N-1} \\
f_N
\end{bmatrix} = \begin{bmatrix}
A \\
(\Delta x)^2 r_1 \\
(\Delta x)^2 r_2 \\
\vdots \\
(\Delta x)^2 r_{N-1} \\
B
\end{bmatrix}$$

We can modify the first and last rows of L in the case that we have different boundary values.

# Chapter 2: Subspaces, Bases, and Dimension

#### 2.1 Vector spaces and subspaces

• If V is a vector space, then  $S \subseteq V$  is a subspace of V if  $\forall u, v \in S$  and a, b scalars,  $au + bv \in S$ .

#### 2.2 The four fundamental subspaces

- The four subspaces
  - $-N(A) \subseteq \mathbb{R}^n$ 
    - \* Solve the homogenous linear system Ax = 0. This is the same as solving Ux = 0, where U = rref(A).
  - $-R(A) \subseteq R^m$ 
    - \* Find the linearly independent columns of A, which correspond to pivot columns of U.
    - \* Use  $R(A) = R(W^T)$ , where  $W = \text{rref}(A^T)$

- $-N(A^T)\subseteq R^m$ 
  - \* Solve the homogenous linear system  $A^T x = 0$ . This is the same as solving W x = 0.
- $-R(A^T) \subseteq R^n$ 
  - \* Find the linearly independent columns of  $A^T$ , which correspond to linearly independent columns of W
  - \* Use  $R(A^T) = R(U^T)$
- Dimensions

$$-\operatorname{rank}(A^T) = \operatorname{rank}(A)$$

$$-\operatorname{rank}(A) + \dim(N(A)) = n$$

$$-\operatorname{rank}(A^T) + \dim(N(A^T)) = m$$

• Properties of the inner product

$$-\langle \mathbf{x}, \mathbf{y} \rangle = \mathbf{x}^T \mathbf{y}$$

$$-\ \langle \mathbf{x}, \mathbf{y} \rangle = \langle \mathbf{y}, \mathbf{x} \rangle$$

$$-\langle a\mathbf{x} + b\mathbf{y}, \mathbf{z} \rangle = a\langle \mathbf{x}, \mathbf{z} \rangle + b\langle \mathbf{y}, \mathbf{z} \rangle$$

$$- \langle \mathbf{x}, \mathbf{x} \rangle = \|x\|_2^2$$

- Two subspaces are orthogonal iff their basis vectors are pairwise orthogonal. If B is a basis for  $S_1$  and C is a basis for  $S_2$ , then  $B^TC = 0$ .
- Relation among  $N(A), N(A^T), R(A), R(A^T)$

$$-N(A) = [R(A^T)]^{\perp} N(A^T) = [R(A)]^{\perp}$$

- In order to determine whether **x** is in R(A), we can check if  $\mathbf{x} \perp N(A^T)$ 

## 2.3 Graphs and networks

- D is the  $m \times n$  incidence matrix where m = # edges, n = # nodes
- $\mathbf{v} \in \mathbb{R}^n$  assigns the voltage to each node.  $D\mathbf{v} \in \mathbb{R}^m$  assigns the voltage difference across each edge.  $\mathbf{y} \in \mathbb{R}^m$  assigns a current to each edge.  $D^T\mathbf{y} \in \mathbb{R}^n$  assigns the current accumulation at each node.
- The four subspaces
  - N(D) is the set of voltages on nodes such that the voltage difference across each edge is 0. Since the voltages of any connected nodes have to be identical,  $N(D) = \begin{bmatrix} 1 \\ \vdots \\ 1 \end{bmatrix}$  for any connected graph.

$$\dim(N(D)) = \#$$
 connected components

-R(D) is the set of possible voltage differences across edges. Since the sum of differences around a closed loop is 0, we have one constraining equation for every linearly independent loop in the graph.

$$\dim(R(D)) = m - \#$$
 independent loops

 $-N(D^T)$  is the set of edge currents that do not result in any current building up at any nodes. We can find a basis by using the fact that the current around a loop will not build up at any node.

$$\dim(N(D^T)) = \#$$
 independent loops

 $-R(D^T)$  is the set of possible current accumulations at each node. Since the current building up at one node must have come from other nodes, in a connected graph, we only have one constraining equation.

$$\dim(R(D^T)) = n - \#$$
 connected components

- The Laplacian:  $L = D^T R^{-1} D$ 
  - L is symmetrical
  - If all resistors are  $1\Omega$ , then

$$L_{ii} =$$
 of neighbors of node  $i$ 

$$L_{ij} = \begin{cases} -1 & \text{if vertex } i, j \text{ are connected} \\ 0 & \text{otherwise} \end{cases}$$

- More generally,

$$L_{ii} = \sum_{j} \frac{1}{R_{j}}$$
 sum over all edges connecting to node i

$$L_{ij} = \begin{cases} -\frac{1}{R_k} & \text{where } R_k \text{ is the resistance between nodes } i, j \\ 0 & \text{otherwise} \end{cases}$$

- -N(L) = N(D)
- Kirchoff's law:  $L\mathbf{v} = 0$  when there are no batteries involved.
- When we connect a battery, we have the system

$$L\mathbf{v} = \mathbf{z}$$

where z represents the current accumulating at each node.

$$\begin{bmatrix} A & B^T \\ B & C \end{bmatrix} \begin{bmatrix} \mathbf{b} \\ \mathbf{v} \end{bmatrix} \begin{bmatrix} \mathbf{J} \\ \mathbf{0} \end{bmatrix}$$

$$\mathbf{J} = (A - B^T C^{-1} B) \mathbf{b}$$

Since J gives the current vector and we know the voltage of the battery, we can solve

$$R_{eq} = \frac{\text{current}}{\text{voltage}}$$

# Chapter 3: Orthogonality

# 3.1 Projections

• Projection of  $\mathbf{y}$  onto L is the vector in L that is closest to  $\mathbf{y}$ 

$$P_{\mathbf{x}}(\mathbf{y}) = \frac{\langle \mathbf{x}, \mathbf{y} \rangle}{\|\mathbf{x}\|^2} \mathbf{x} = \langle \hat{\mathbf{x}}, \mathbf{y} \rangle \,\hat{\mathbf{x}}$$

$$R(P) = L \quad N(P) = L^{\perp}$$

• Q = I - P projects onto  $L^{\perp}$ 

• P is an orthogonal projection matrix if  $P^2 = P$  and  $P^T = P$ 

$$P_{\mathbf{x}} = \frac{\mathbf{x}\mathbf{x}^T}{\left\|\mathbf{x}\right\|^2}$$

$$P_{R(A)} = A(A^T A)^{-1} A^T = \tilde{A} (\tilde{A}^T \tilde{A})^{-1} \tilde{A}^T$$

• Least squares equation to solve  $A\mathbf{x} = \mathbf{b}$  when there is no solution

$$A^T A \mathbf{x}^* = A^T \mathbf{b}$$

 $-A^TA$  is invertible iff A is invertible

$$\mathbf{x}_{LS} = (A^T A)^{-1} A^T \mathbf{b}$$

- Pseudoinverse of A:  $(A^T A)^{-1} A^T$
- Polynomial regression m < n

$$p(x) = a_1 x^{n-1} + \dots + a_{n-1} x + a_n$$

$$V\mathbf{a} = \mathbf{y}$$

$$V = \begin{bmatrix} x_1^{m-1} & \dots & x_1 & 1 \\ \vdots & & \vdots & \vdots \\ x_n^{m-1} & \dots & x_n & 1 \end{bmatrix} \quad \mathbf{a} = \begin{bmatrix} a_1 \\ \vdots \\ a_m \end{bmatrix} \quad \mathbf{y} = \begin{bmatrix} y_1 \\ \vdots \\ y_n \end{bmatrix}$$

$$\mathbf{a}_{LS} = (V^T V)^{-1} V^T \mathbf{y}$$

# 3.2 Complex vector spaces and inner product

- Complex number z = a + bi
  - Modulus:  $|z| = \sqrt{a^2 + b^2}$
  - Complex conjugate:  $\overline{z} = a bi$
  - Argument:  $\theta = \arctan\left(\frac{b}{a}\right)$
- Some formulas

$$-z = |z|(\cos\theta + i\sin\theta) = Re^{i\theta}$$

$$-e^{i\theta} = \cos\theta + i\sin\theta$$

• Complex inner product

$$\langle z,z\rangle=\overline{z_1}\cdot z_1+\ldots+\overline{z_n}\cdot z_n$$

• Properties

$$-\langle s\mathbf{v}, \mathbf{w} \rangle = \overline{s} \langle \mathbf{v}, \mathbf{w} \rangle$$

$$-\langle \mathbf{v}, s\mathbf{w} \rangle = s\langle \mathbf{v}, \mathbf{w} \rangle$$

$$-\langle \mathbf{v}, \mathbf{w} \rangle = \overline{\langle \mathbf{w}, \mathbf{v} \rangle}$$

$$-\langle \mathbf{v}, \mathbf{w} \rangle = \overline{\mathbf{v}}^T \mathbf{w}$$

$$-\langle A\mathbf{v}, \mathbf{w} \rangle = \langle \mathbf{v}, A^*\mathbf{w} \rangle$$

• Adjoint of A:  $A^* = \overline{A}^T$ 

### 3.3 Orthonormal bases, orthogonal & unitary matrices

 $\bullet$  A basis  $\{{\bf q}_1,...,{\bf q}_n\}$  is an orthonormal basis for V if

$$\langle \mathbf{q}_i, \mathbf{q}_j \rangle = \begin{cases} 0 & i \neq j \\ 1 & i = j \end{cases}$$

• For any  $\mathbf{v} \in V$ , there exists a unique set of scalars  $c_1, ..., c_n$  such that

$$\mathbf{v} = c_1 \mathbf{q}_1 + \dots + c_n \mathbf{q}_n = \sum_{i=1}^n \left\langle \mathbf{q}_j, \mathbf{v} \right\rangle \mathbf{q}_j$$

- $\bullet$  A square matrix Q whose columns form an ONB is called
  - Orthogonal if all entries are real
  - Unitary if entries are complex
- Properties
  - $-QQ^* = Q^*Q = I$
  - Rows of Q is also an orthogonal/unitary matrix
  - $-\left\| Q\mathbf{v}\right\| _{2}=\left\| \mathbf{v}\right\| _{2}$

#### 3.4 Fourier series

• 1-period signals are modeled by

$$L^{2}[0,1] = \left\{ f: [0,1] \to C: \int_{0}^{1} |f(t)|^{2} dt < \infty \right\}$$

• The inner product

$$\langle f, g \rangle = \int_0^1 \overline{f(t)} g(t) dx$$
$$\langle f, f \rangle = \int_0^1 |f(x)|^2 dx = ||f||_2^2$$

• An orthonormal basis for  $L^2[0,1]$  is given by

$$\{e^{2\pi int}:n\in Z\}$$

• So any  $f \in L^2[0,1]$  can be decomposed as

$$f(t) = \sum_{n = -\infty}^{\infty} c_n e^{2\pi i nt}$$

$$c_n = \left\langle e^{2\pi i n t}, f \right\rangle = \int_0^1 f(t) e^{-2\pi i n t} dt$$

• The average value of f(x) on [0,1] is given by

$$\int_{0}^{1} f(t)dt = \int_{0}^{1} f(t)e^{-2\pi i \cdot 0 \cdot t} = c_{0}$$

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•  $c_k$  and  $c_{-k}$  correspond to the frequency content of f(x) at frequency  $2\pi |k|$ 

• If f(t) is periodic with period T, then we have  $L^2[0,T]$  and an ONB is given by

$$\left\{ \frac{1}{\sqrt{T}} e^{2\pi i n t/T} : n \in Z \right\}$$

$$f(t) = \sum_{n=-\infty}^{\infty} c_n e^{2\pi i n t/T} dt$$

$$c_n = \frac{1}{T} \left\langle e^{2\pi i n t/T}, f \right\rangle = \frac{1}{T} \int_0^T f(t) e^{-2\pi i n t/T} dt$$

• Parseval's formula

$$\sum_{n=-\infty}^{\infty} |c_n|^2 = \frac{1}{T} \int_0^T |f(t)|^2 dt$$
$$\sum_{n=-\infty}^{\infty} |c_n|^2 = \int_0^1 |f(t)|^2 dt$$

## Chapter 4: Eigenvalues and Eigenvectors

### 4.1 Eigenvalues and eigenvectors

- Eigenvalue/eigenvector pair of  $n \times m$  matrix:  $A\mathbf{v} = \lambda \mathbf{v}$
- Eigenvalues are the roots of the characteristic polynomial  $p(\lambda) = \det(A \lambda I)$ . For each  $\lambda_j$ , solve  $(A \lambda_j I)\mathbf{v} = \mathbf{0}$ . This subspace is the eigenspace  $E_{\lambda_j}$ .
- The algebraic multiplicity  $m_j$  of an eigenvalue  $\lambda_j$  is the power to which  $(\lambda \lambda_j)$  divides the polynomial. The geometric multiplicity  $d_j$  is  $\dim(E_{\lambda_j})$ .  $1 \le d_j \le m_j$  for all j.
- An eigenbasis exists iff  $d_j = m_j$  for each j. In particular, if there are n distinct eigenvalues, then  $d_j = m_j = 1$  for each j.
- A is diagonalizable if it has an eigenbasis

$$A = SDS^{-1}$$

• From diagonalization, we easily see that

$$\det(A)=\lambda_1\cdots\lambda_n$$
 
$$\operatorname{tr}(A)=\lambda_1+\ldots+\lambda_n$$
 
$$A^k=SD^kS^{-1} \text{ for all } k\in Z^+ \text{ and for } k\in Z^- \text{ if } D^{-1} \text{ exists}$$

#### 4.2 Hermitian matrices

• A square matrix is Hermitian if

$$A^* = A$$

A Hermitian matrix that is real is also symmetric.

- A is Hermitian iff  $\langle A\mathbf{v}, \mathbf{w} \rangle = \langle \mathbf{v}, A\mathbf{w} \rangle$
- A Hermitian matrix has
  - Real eigenvalues
  - Eigenvectors that correspond to distinct eigenvalues are orthogonal
- Every Hermitian matrix is unitarily diagonalizable

$$A = UDU^*$$

#### 4.3 Power method

• Assumptions

$$- |\lambda_1| > |\lambda_2| \ge \dots \ge |\lambda_n|$$

$$-\{\mathbf{v}_1,...,\mathbf{v}_n\}$$
 is an eigenbasis

$$-\|\mathbf{v}_j\|=1$$
 for all  $j$ 

• Algorithm

$$\mathbf{x}_0 = c_1 \mathbf{v}_1 + \dots + c_n \mathbf{v}_n \quad \text{random}$$

$$A^k \mathbf{x}_0 \approx \lambda_1^k c_1 \mathbf{v}_1$$

$$\mathbf{v}_1 = \frac{A^k \mathbf{x}_0}{\|A^k \mathbf{x}_0\|}$$

$$\lambda_1 = \langle \mathbf{v}_1, A \mathbf{v}_1 \rangle$$

- ullet The random vector  $\mathbf{x}_0$  should have a component in the  $\mathbf{v}_1$  direction, otherwise power method fails
- To find the eigenvalue closest to s, use the power method on  $(A sI)^{-1}$ . This returns  $\mathbf{v}_{j^*}$ , and we find  $\lambda_{j^*} = \langle \mathbf{v}_{j^*}, A\mathbf{v}_{j^*} \rangle$ .

### 4.4 Recursion relations

• Consider the Fibonacci sequence:  $F_{n+1} = F_n + F_{n-1}$ 

$$\begin{bmatrix} F_{n+1} \\ F_n \end{bmatrix} = \begin{bmatrix} 1 & 1 \\ 1 & 0 \end{bmatrix}^k \begin{bmatrix} 1 \\ 0 \end{bmatrix}$$

- Then diagonalize A to easily compute  $A^k$
- If there exists  $\lambda_i = e^{i\frac{2\pi}{T}}$ , then choosing  $\mathbf{x}_0 = \mathbf{v}_i$  causes periodic oscillations of period T

#### 4.6 Markov chains

- $p_{ij}$  is the transition probability from j to i
- $x_{n,i}$  is the probability of being at state i at time n. Define the state vector  $\mathbf{x}_n = \begin{bmatrix} x_{n,1} \\ x_{n,2} \\ x_{n,3} \end{bmatrix}$
- $\bullet$  After *n* iterations,

$$\mathbf{x}_n = P^n \mathbf{x}_0$$

• A matrix P is stochastic if

$$P = \begin{bmatrix} p_{11} & \dots & p_{1k} \\ \vdots & & \vdots \\ p_{k1} & \dots & p_{kk} \end{bmatrix}$$

such that

$$0 \le p_{ij} \le 1$$
 (non-zero probabilities)

$$\sum_{i=1}^{k} p_{ij} = 1 \text{ for each } j \pmod{\text{sum is 1}}$$

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• When n is large and  $\lim_{n\to\infty} P^n \mathbf{x}$  exists,

$$P^n \mathbf{x}_0 = P\mathbf{x} = \mathbf{x}$$

So  $P\mathbf{x} = \mathbf{x}$  and  $\mathbf{x}$  is a stationary point and corresponds to  $\lambda = 1$ . Eigenvector  $\mathbf{v}_1$  corresponding to  $\lambda_1 = 1$  has non-negative entries. All other eigenvalues of P satisfy  $|\lambda_j| \leq 1$ . Eigenvectors corresponding to  $|\lambda_j| < 1$  have entries that add up to 0.

- If P or  $P^k$  for some  $k \in N$  has all positive entries, then  $\lambda_1$  is the dominant eigenvalue, and normalizing  $\mathbf{v}_1$  so that its 1-norm is 1 gives the stationary point. This occurs when there is a nonzero probability of ending up anywhere in the graph.
- For Google pagerank, we define a stochastic matrix that represents the probability of each page linking to every other page. If there are no links out of a node, then we assign equal probabilities from that node to any node (including staying at itself).
- Because P often has so many 0 entries, it almost always has many eigenvalues  $|\lambda| = 1$ , so we can't use power method to compute  $\mathbf{v}_1$ . To remedy this, we choose the damping factor  $0 \le \alpha \le 1$  and use the Google matrix.

$$G = \alpha S + (1 - \alpha)Q \quad \text{where} \quad Q = \begin{bmatrix} \frac{1}{N} & \dots & \frac{1}{N} \\ \vdots & & \vdots \\ \frac{1}{N} & \dots & \frac{1}{N} \end{bmatrix}$$

This guarantees that  $|\lambda_1| = 1$  is the dominant eigenvalue and the power method can be used.

### 4.7 Singular value decomposition

• Singular value decomposition generalizes diagonalization to non-square matrices.

$$A = U \cdot \Sigma \cdot V^*$$

- Observations that help us derive the SVD
  - All eigenvalues of  $A^*A$  are non-negative
  - $-A^*A$  and  $AA^*$  have identical non-zero eigenvalues (but may have different numbers of zero eigenvalues)
  - Both  $A^*A$  and  $AA^*$  are Hermitian (and thus unitarily diagonalizable)
- We get  $A = U\Sigma V^*$ , where
  - V is formed by eigenvectors of A\*A

$$A^*A = V\Sigma_1^2 V^*$$

- U is formed by eigenvectors of  $AA^*$ 

$$AA^* = U\Sigma_2^2 U^*$$

-  $\Sigma$  is formed by the min $\{m,n\}$  eigenvalues of  $A^*A$  and  $AA^*$ 

$$\Sigma = \begin{bmatrix} \sigma_1 & & & 0 & \dots & 0 \\ & \sigma_2 & & 0 & \dots & 0 \\ & & \ddots & & \vdots & & \vdots \\ & & \sigma_m & 0 & \dots & 0 \end{bmatrix} \quad \text{for } m < n$$

 $\bullet$  Let r be the number of non-zero singular values of A

$$-N(A) = {\mathbf{v}_{r+1}, ..., \mathbf{v}_n}$$

$$- R(A) = {\mathbf{u}_1, ..., \mathbf{u}_r} - N(A^*) = {\mathbf{u}_1, ..., \mathbf{u}_r}^{\perp} = {\mathbf{u}_{r+1}, ..., \mathbf{u}_m}$$

• Matrix norms

$$- ||A|| = \sigma_1$$
  
-  $||A||_{HS}^2 = \sigma_1^2 + \dots + \sigma_r^2$ 

 $\bullet$  The left-inverse of A is given by

$$A^\dagger = V \Sigma^\dagger U^*$$

 $\quad \text{where} \quad$ 

$$\Sigma^{\dagger} = \begin{bmatrix} \frac{1}{\sigma_1} & & & \\ & \ddots & & \\ & & \frac{1}{\sigma_n} \\ 0 & \dots & 0 \\ \vdots & & \vdots \\ 0 & \dots & 0 \end{bmatrix}$$