# Analysis review: O notation, Taylor series, and linear algebra 

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## Introduction

One final lecture of analysis review, in which we go over three indispensable tools that we will use constantly in the remainder of the course:

- $O$, o notation
- Taylor series expansions
- Linear algebra


## o-notation: Motivation

- When investigating the asymptotic behavior of functions, it is often convenient to replace unwieldy expressions with compact notation
- For example, if we encountered the mathematical expression

$$
x^{2}+a-a
$$

we would obviously want to replace it with $x^{2}$ since $a-a=0$

- However, what if we encounter something like

$$
x^{2}+\frac{5 \theta}{\sqrt{n}}-\frac{3 \theta}{n+5} ?
$$

- We can no longer just replace this with $x^{2}$


## o-notation: Motivation (cont'd)

- However, as $n$ gets larger, the expression gets closer and closer to $x^{2}$
- It would be convenient to have a shorthand notation for this, something like $x^{2}+o_{n}$, where $o_{n}$ represents some quantity that becomes negligible as $n$ becomes large
- This is the basic idea behind $o$-notation, and its simplifying powers become more apparent as the mathematical expression we are dealing with becomes more complicated:

$$
\frac{x^{2}+\frac{5 \theta}{\sqrt{n}}-\frac{3 \theta}{n+5}}{\left(n^{2}+5 n-2\right) /\left(n^{2}-3 n+1\right)}+\frac{\exp \left\{-\frac{1}{2}\|\mathbf{x}-\boldsymbol{\mu}\|^{2}\right\}}{2 \sqrt{n} \theta \int_{0}^{\infty} g(s) d s}
$$

## o-notation

- There is where something called o-notation comes in: a formal way of handling terms that effectively "cancel out" as we take limits
- Definition: A sequence of numbers $x_{n}$ is said to be $o(1)$ if it converges to zero. Likewise, $x_{n}$ is said to be $o\left(r_{n}\right)$ if

$$
\frac{x_{n}}{r_{n}} \rightarrow 0
$$

as $n \rightarrow \infty$.

- When the rate is constant, o notation is pretty straightforward:

$$
x^{2}+\frac{5 \theta}{\sqrt{n}}-\frac{3 \theta}{n+5}=x^{2}+o(1)
$$

## o-notation remarks

- When the rate is not constant, expressions are a bit harder to think about - it helps to go over some cases:
- For example:
- $x_{n} \rightarrow \infty$, but $r_{n} \rightarrow \infty$ even faster:

$$
n=o\left(n^{2}\right)
$$

- $r_{n} \rightarrow 0$, but $x_{n} \rightarrow 0$ even faster:

$$
\frac{1}{n^{2}}=o(1 / n)
$$

## $O$-notation

- A very useful companion of $o$-notation is $O$-notation, which denotes whether or not a term remains bounded as $n \rightarrow \infty$
- Definition: A sequence of numbers $x_{n}$ is said to be $O(1)$ if there exist $M$ and $N$ such that

$$
\left|x_{n}\right|<M
$$

for all $n>N$. Likewise, $x_{n}$ is said to be $O\left(r_{n}\right)$ if there exist $M$ and $N$ such that for all $n>N$,

$$
\left|\frac{x_{n}}{r_{n}}\right|<M .
$$

## O-notation remarks

- For example,

$$
\frac{\exp \left\{-\frac{1}{2}\|\mathbf{x}-\boldsymbol{\mu}\|^{2}\right\}}{2 \sqrt{n} \theta \int_{0}^{\infty} g(s) d s}=O\left(n^{-1 / 2}\right)
$$

- Note that $x_{n}=O(1)$ does not necessarily mean that $x_{n}$ is bounded, just that it is eventually bounded
- Note also that just because a term is $O(1)$, this does not necessarily mean that it has a limit; for example,

$$
\sin \left(\frac{n \pi}{2}\right)=O(1)
$$

even though the sequence does not converge

## O-notation remarks (cont'd)

- You may encounter the ambiguous phrase " $x_{n}$ is of order $r_{n}$ "
- The author may mean that $x_{n}=O\left(r_{n}\right)$
- However, it might also mean something stronger: that there exist positive constants $m$ and $M$ such that

$$
m \leq\left|\frac{x_{n}}{r_{n}}\right| \leq M
$$

for large enough $n$; i.e., the ratio is bounded above but also bounded below

- In other words, $x_{n}=O\left(r_{n}\right)$ but in addition $x_{n} \neq o\left(r_{n}\right)$; some authors use the notation $x_{n} \asymp r_{n}$ to denote this situation


## Informative-ness of $o$ and $O$ notation

- There are typically many ways of writing an expression using $O$ notation, although not all of them will be equally informative
- For example, if $x_{n}=1+\frac{1}{n}$, then all of the following are true:

$$
\begin{aligned}
& x_{n}=1+o(1) \\
& x_{n}=1+O(1) \\
& x_{n}=1+O\left(\frac{1}{n}\right) \\
& x_{n} \asymp 1+O\left(\frac{1}{n}\right)
\end{aligned}
$$

(least informative)
(more informative)
(most informative)

## Algebra of $O, o$ notation

$O, o$-notation are useful in combination because simple rules govern how they interact with each other
Theorem: For $a \leq b$ :

$$
\begin{aligned}
O(1)+O(1) & =O(1) & O\{O(1)\} & =O(1) \\
o(1)+o(1) & =o(1) & o\{O(1)\} & =o(1) \\
o(1)+O(1) & =O(1) & o\left(r_{n}\right) & =r_{n} o(1) \\
O(1) O(1) & =O(1) & O\left(r_{n}\right) & =r_{n} O(1) \\
O(1) o(1) & =o(1) & O\left(n^{a}\right)+O\left(n^{b}\right) & =O\left(n^{b}\right) \\
\{1+o(1)\}^{-1} & =O(1) & o\left(n^{a}\right)+o\left(n^{b}\right) & =o\left(n^{b}\right)
\end{aligned}
$$

## Remarks

- $O, o$ "equations" are meant to be read left-to-right; for example, $O(\sqrt{n})=O(n)$ is a valid statement, but $O(n)=O(\sqrt{n})$ is not
- Exercise: Determine the order of

$$
n^{-2}\left\{(-1)^{n} \sqrt[n]{2}+n\left(1+\frac{1}{n}\right)^{n}\right\}
$$

- As we will see in a week or two, there are stochastic equivalents of these concepts, involving convergence in probability and being bounded in probability
- As such, we won't do a great deal with $O, o$-notation right now, but will use the stochastic equivalents extensively


## Taylor series: Introduction

As we will see (many times!), it is useful to be able to approximate a complicated function with a simple polynomial (this is the idea behind Taylor series approximation):


## Taylor series: Introduction (cont'd)

- It is difficult to overstate the importance of Taylor series expansions to statistical theory, and for that reason we are now going to cover them fairly extensively
- In particular, Taylor's theorem comes in a number of versions, and it is worth knowing several of them, since they come up in statistics quite often
- Furthermore, students often have not seen the multivariate versions of these expansions


## Taylor's theorem

- Theorem (Taylor): Suppose $n$ is a positive integer and $f: \mathbb{R} \rightarrow \mathbb{R}$ is $n$ times differentiable at a point $x_{0}$. Then

$$
f(x)=\sum_{k=0}^{n} \frac{f^{(k)}\left(x_{0}\right)}{k!}\left(x-x_{0}\right)^{k}+R_{n}\left(x, x_{0}\right)
$$

where the remainder $R_{n}$ satisfies

$$
R_{n}\left(x, x_{0}\right)=o\left(\left|x-x_{0}\right|^{n}\right) \text { as } x \rightarrow x_{0}
$$

- If $f^{(n+1)}\left(x_{0}\right)$ exists, you could also say that $R_{n}$ is $O\left(\left|x-x_{0}\right|^{n+1}\right)$
- This form of the remainder is sometimes called the Peano form


## Taylor's theorem: Lagrange form

- Theorem (Taylor): Suppose $f: \mathbb{R} \rightarrow \mathbb{R}$ is $n+1$ times differentiable on an open interval containing $x_{0}$. Then for any point $x$ in that interval, there exists $\bar{x} \in\left(x, x_{0}\right)$ :

$$
R_{n}\left(x, x_{0}\right)=\frac{f^{(n+1)}(\bar{x})}{(n+1)!}\left(x-x_{0}\right)^{n+1}
$$

- This is also known as the mean-value form, as the mean value theorem is the central idea in proving the result


## Comparing the two forms

- Comparing the Basic and Lagrange forms for a second-order expansion,

$$
\begin{aligned}
& f\left(x_{0}\right)+f^{\prime}\left(x_{0}\right)\left(x-x_{0}\right)+\frac{1}{2} f^{\prime \prime}\left(x_{0}\right)\left(x-x_{0}\right)^{2}+o\left(\left|x-x_{0}\right|^{2}\right) \\
& f\left(x_{0}\right)+f^{\prime}\left(x_{0}\right)\left(x-x_{0}\right)+\frac{1}{2} f^{\prime \prime}(\bar{x})\left(x-x_{0}\right)^{2}
\end{aligned}
$$

- We can see that in the second case, we have a simpler expression, but to obtain it, we require $f^{\prime \prime}$ to exist along the entire interval from $x$ to $x_{0}$, not just at the point $x_{0}$


## Example: Absolute value

- For example, consider approximating the function $f(x)=|x|$ at $x_{0}=-0.1$
- Note that $f^{\prime}$ exists at $x_{0}$, but not at 0
- The basic form of Taylor's theorem says that if we get close enough to $x_{0}$, the approximation $f(-0.1)+f^{\prime}(-0.1)(x+0.1)$ becomes very accurate - indeed, the remainder is exactly zero for any $x$ within 0.1 of $x_{0}$
- However, suppose $x=0.2$; since $f$ is not differentiable at zero, we are not guaranteed the existence of a point $\bar{x}$ such that

$$
f(0.2)=f(-0.1)+0.3 f^{\prime}(\bar{x}) ;
$$

and indeed in this case no such point exists

## Lagrange bound

- One reason why the Lagrange form is more powerful is that it allows us to establish error bounds - to know exactly how close $x$ must be to $x_{0}$ in order to ensure that the approximation error is less than $\epsilon$
- In particular, if there exists an $M$ such that $\left|f^{(n+1)}(x)\right| \leq M$ over the interval $\left(x, x_{0}\right)$, then

$$
\left|R_{n}(x)\right| \leq \frac{M}{(n+1)!}\left|x-x_{0}\right|^{n+1}
$$

## Multivariable forms of Taylor's theorem

- We now turn our attention to the multivariate case
- For the sake of clarity, I'll present the first- and second-order expansions for each of the previous forms, rather than abstract formulae involving $f^{(n)}$
- Lastly, I'll provide a form that goes out to third order, although higher orders are less convenient as they can't be represented compactly using vectors and matrices
- Note that these forms are only covering the case of scalar-valued functions $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$; we will need results for the vector-valued case $f: \mathbb{R}^{d} \rightarrow \mathbb{R}^{k}$ as well, but we will go over that in a later lecture


## Taylor's theorem

- Theorem (Taylor): Suppose $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ is differentiable at a point $\mathbf{x}_{0}$. Then

$$
f(\mathbf{x})=f\left(\mathbf{x}_{0}\right)+\nabla f\left(\mathbf{x}_{0}\right)^{\top}\left(\mathbf{x}-\mathbf{x}_{0}\right)+o\left(\left\|\mathbf{x}-\mathbf{x}_{0}\right\|\right)
$$

- Theorem (Taylor): Suppose $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ is twice differentiable at a point $\mathbf{x}_{0}$. Then

$$
\begin{aligned}
f(\mathbf{x})= & f\left(\mathbf{x}_{0}\right)+\nabla f\left(\mathbf{x}_{0}\right)^{\top}\left(\mathbf{x}-\mathbf{x}_{0}\right)+ \\
& \frac{1}{2}\left(\mathbf{x}-\mathbf{x}_{0}\right)^{\top} \nabla^{2} f\left(\mathbf{x}_{0}\right)\left(\mathbf{x}-\mathbf{x}_{0}\right)+o\left(\left\|\mathbf{x}-\mathbf{x}_{0}\right\|^{2}\right)
\end{aligned}
$$

## Taylor's theorem: Lagrange form

- Theorem (Taylor): Suppose $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ is differentiable on $N_{r}\left(\mathbf{x}_{0}\right)$. Then for any $\mathbf{x} \in N_{r}\left(\mathbf{x}_{0}\right)$, there exists $\overline{\mathrm{x}}$ on the line segment connecting x and $\mathrm{x}_{0}$ such that

$$
f(\mathbf{x})=f\left(\mathbf{x}_{0}\right)+\nabla f(\overline{\mathbf{x}})^{\top}\left(\mathbf{x}-\mathbf{x}_{0}\right)
$$

- Theorem (Taylor): Suppose $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ is twice differentiable on $N_{r}\left(\mathbf{x}_{0}\right)$. Then for any $\mathbf{x} \in N_{r}\left(\mathbf{x}_{0}\right)$, there exists $\overline{\mathbf{x}}$ on the line segment connecting $\mathbf{x}$ and $\mathbf{x}_{0}$ such that

$$
\begin{aligned}
f(\mathbf{x})= & f\left(\mathbf{x}_{0}\right)+\nabla f\left(\mathbf{x}_{0}\right)^{\top}\left(\mathbf{x}-\mathbf{x}_{0}\right)+ \\
& \frac{1}{2}\left(\mathbf{x}-\mathbf{x}_{0}\right)^{\top} \nabla^{2} f(\overline{\mathbf{x}})\left(\mathbf{x}-\mathbf{x}_{0}\right)
\end{aligned}
$$

- " $\overline{\mathrm{x}}$ on the line segment connecting x and $\mathrm{x}_{0}$ " means that there exists $w \in[0,1]$ such that $\overline{\mathbf{x}}=w \mathbf{x}+(1-w) \mathbf{x}_{0}$


## Taylor's theorem: Third order

Theorem (Taylor): Suppose $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ is three times differentiable on $N_{r}\left(\mathbf{x}_{0}\right)$. Then for any $\mathbf{x} \in N_{r}\left(\mathbf{x}_{0}\right)$, there exists $\overline{\mathbf{x}}$ on the line segment connecting $\mathbf{x}$ and $\mathbf{x}_{0}$ such that

$$
\begin{aligned}
f(\mathbf{x})= & f\left(\mathbf{x}_{0}\right)+\sum_{j=1}^{d} \frac{\partial f\left(\mathbf{x}_{0}\right)}{\partial x_{j}}\left(x_{j}-x_{0 j}\right) \\
& +\frac{1}{2} \sum_{j=1}^{d} \sum_{k=1}^{d} \frac{\partial^{2} f\left(\mathbf{x}_{0}\right)}{\partial x_{j} \partial x_{k}}\left(x_{j}-x_{0 j}\right)\left(x_{k}-x_{0 k}\right) \\
& +\frac{1}{6} \sum_{j=1}^{d} \sum_{k=1}^{d} \sum_{\ell=1}^{d} \frac{\partial^{3} f(\overline{\mathbf{x}})}{\partial x_{j} \partial x_{k} \partial x_{\ell}}\left(x_{j}-x_{0 j}\right)\left(x_{k}-x_{0 k}\right)\left(x_{\ell}-x_{0 \ell}\right),
\end{aligned}
$$

where $\partial f\left(\mathbf{x}_{0}\right) / \partial x_{j}$ is shorthand for $\partial f(\mathbf{x}) / \partial x_{j}$ evaluated at $\mathbf{x}_{0}$

## Linear algebra

- Our last mathematical topic to review/reference is linear algebra, which we will use right away in our next lecture on the multivariate normal distribution
- If this material is new to you, then I would encourage you to read this as well
- It is often useful to switch the rows and columns of a matrix around. The resulting matrix is called the transpose of the original matrix, and denoted with a superscript ${ }^{\top}$ or an apostrophe ':

$$
\mathbf{M}=\left[\begin{array}{rr}
3 & 2 \\
4 & -1 \\
-1 & 2
\end{array}\right] \quad \mathbf{M}^{\top}=\left[\begin{array}{rrr}
3 & 4 & -1 \\
2 & -1 & 2
\end{array}\right]
$$

## Inner and outer products

- Two particularly common situations in which transposes arise are when we multiply a vector by itself
- There are two ways to do this:
- Inner product: $\mathbf{x}^{\top} \mathbf{x}$; note that this is a scalar:

$$
\mathbf{x}^{\top} \mathbf{x}=\sum_{j} x_{j}^{2}
$$

- Outer product: $\mathbf{x x}^{\top}$; note that this is a $d \times d$ matrix (where $d$ is the dimension of $\mathbf{x}$ ):

$$
\left(\mathbf{x} \mathbf{x}^{\top}\right)_{i j}=x_{i} x_{j}
$$

- Note that $\mathrm{x}^{2}$ has no meaning; never write this


## Linear and quadratic forms

- Matrix products involving linear and quadratic forms come up very often in statistics, and it is important to have an intuitive grasp on what they represent
- Here are some useful identities/relationships involving matrix products and their scalar representations:

$$
\begin{aligned}
& \mathbf{a}^{\top} \mathbf{x}=\sum_{i} a_{i} x_{i} ; \quad \mathbf{1}^{\top} \mathbf{x}=\sum_{i} x_{i} \\
& \mathbf{A}^{\top} \mathbf{x}=\left(\begin{array}{lll}
\sum_{i} a_{i 1} x_{i} & \cdots & \sum_{i} a_{i k} x_{i}
\end{array}\right)^{\top} \\
& \mathbf{a}^{\top} \mathbf{W} \mathbf{x}=\sum_{i} \sum_{j} a_{i} w_{i j} x_{j} ; \quad \mathbf{a}^{\top} \mathbf{1} \mathbf{x}=\sum_{i} \sum_{j} a_{i} x_{j} \\
& (\mathbf{A W B})_{i j}=\sum_{k} \sum_{m} a_{i k} w_{k m} b_{m j}
\end{aligned}
$$

## Inverses

- Definition: The inverse of an $n \times n$ matrix $\mathbf{A}$, denoted $\mathbf{A}^{-1}$, is the matrix satisfying $\mathbf{A A}^{-1}=\mathbf{A}^{-1} \mathbf{A}=\mathbf{I}_{n}$, where $\mathbf{I}_{n}$ is the $n \times n$ identity matrix.
- Note: We're sort of getting ahead of ourselves by saying that $\mathbf{A}^{-1}$ is "the" matrix satisfying $\mathbf{A} \mathbf{A}^{-1}=\mathbf{I}_{n}$, but it is indeed the case that if a matrix has an inverse, the inverse is unique
- Some useful results:

$$
\begin{aligned}
(\mathbf{A}+\mathbf{B})^{\top} & =\mathbf{A}^{\top}+\mathbf{B}^{\top} \\
(\mathbf{A B})^{\top} & =\mathbf{B}^{\top} \mathbf{A}^{\top} \\
(\mathbf{A B})^{-1} & =\mathbf{B}^{-1} \mathbf{A}^{-1} \\
\left(\mathbf{A}^{\top}\right)^{-1} & =\left(\mathbf{A}^{-1}\right)^{\top}
\end{aligned}
$$

## Singular matrices

- However, not all matrices have inverses; for example

$$
\mathbf{A}=\left[\begin{array}{ll}
1 & 2 \\
2 & 4
\end{array}\right]
$$

- There does not exist a matrix such that $\mathbf{A A}^{-1}=\mathbf{I}_{2}$
- Such matrices are said to be singular
- Remark: Only square matrices have inverses; an $n \times m$ matrix A might, however, have a left inverse (satisfying $\mathbf{B A}=\mathbf{I}_{m}$ ) or right inverse (satisfying $\mathbf{A B}=\mathbf{I}_{n}$ )


## Positive definite

- A related notion is that of a "positive definite" matrix, which applies to symmetric matrices
- Definition: A symmetric $n \times n$ matrix $\mathbf{A}$ is said to be positive definite if for all $\mathbf{x} \in \mathbb{R}^{n}$,

$$
\mathbf{x}^{\top} \mathbf{A} \mathbf{x}>0 \quad \text { if } \mathbf{x} \neq 0
$$

- The two notions are related in the sense that if $\mathbf{A}$ is positive definite, then (a) $\mathbf{A}$ is not singular and (b) $\mathbf{A}^{-1}$ is also positive definite
- If $\mathbf{x}^{\top} \mathbf{A} \mathbf{x} \geq 0$, then $\mathbf{A}$ is said to be positive semidefinite
- In statistics, these classifications are particularly important for variance-covariance matrices, which are always positive semidefinite (and positive definite, if they aren't singular)


## Square root of a matrix

- These concepts are important with respect to knowing whether a matrix has a "square root"
- Definition: An $n \times n$ matrix $\mathbf{A}$ is said to have a square root if there exists a matrix $\mathbf{B}$ such that $\mathbf{B B}=\mathbf{A}$.
- Theorem: Let $\mathbf{A}$ be a positive definite matrix. Then there exists a unique matrix $\mathbf{A}^{1 / 2}$ such that $\mathbf{A}^{1 / 2} \mathbf{A}^{1 / 2}=\mathbf{A}$.
- Positive semidefinite matrices have square roots as well, although they aren't necessarily unique


## Rank

- We also need to be familiar with the concept of matrix rank (there are many ways of defining rank; all are equivalent)
- Definition: The rank of a matrix is the dimension of its largest nonsingular submatrix.
- For example, the following $3 \times 3$ matrix is singular, but contains a nonsingular $2 \times 2$ submatrix, so its rank is 2 :

$$
\mathbf{A}=\left[\begin{array}{lll}
1 & 2 & \not 2 \\
\not 2 & A & \not 6 \\
1 & 0 & \not 1
\end{array}\right]
$$

- Note that a nonsingular $n \times n$ matrix has rank $n$, and is said to be full rank


## Rank and multiplication

- There are many results and theorems involving rank; we're not going to cover them all, but it is important to know that rank cannot be increased through the process of multiplication
- Theorem: For any matrices $\mathbf{A}$ and $\mathbf{B}$ with appropriate dimensions, $\operatorname{rank}(\mathbf{A B}) \leq \operatorname{rank}(\mathbf{A})$ and $\operatorname{rank}(\mathbf{A B}) \leq \operatorname{rank}(\mathbf{B})$.
- In particular, $\operatorname{rank}\left(\mathbf{A}^{\top} \mathbf{A}\right)=\operatorname{rank}\left(\mathbf{A} \mathbf{A}^{\top}\right)=\operatorname{rank}(\mathbf{A})$


## Expectation and variance

- In addition, we need some results on expected values of vectors and functions of vectors
- First of all, we need to define expectation and variance as they pertain to random vectors
- Definition: Let $\mathbf{x}=\left(X_{1} X_{2} \cdots X_{d}\right)^{\top}$ denote a vector of random variables, then $\mathbb{E}(\mathbf{x})=\left(\mathbb{E} X_{1} \mathbb{E} X_{2} \cdots \mathbb{E} X_{d}\right)^{\top}$. Meanwhile, $\mathbb{V} \mathbf{x}$ is a $d \times d$ matrix:

$$
\begin{aligned}
\mathbb{V} \mathbf{x} & =\mathbb{E}\left\{(\mathbf{x}-\boldsymbol{\mu})(\mathbf{x}-\boldsymbol{\mu})^{\top}\right\} \text { with elements } \\
(\mathbb{V} \mathbf{x})_{i j} & =\mathbb{E}\left\{\left(X_{i}-\mu_{i}\right)\left(X_{j}-\mu_{j}\right)\right\},
\end{aligned}
$$

where $\mu_{i}=\mathbb{E} X_{i}$. The matrix $\mathbb{V} \mathbf{x}$ is referred to as the variance-covariance matrix of $\mathbf{x}$.

## Linear and quadratic forms

- Letting A denote a matrix of constants and $\mathbf{x}$ a random vector with mean $\boldsymbol{\mu}$ and variance $\boldsymbol{\Sigma}$,

$$
\begin{aligned}
\mathbb{E}\left(\mathbf{A}^{\top} \mathbf{x}\right) & =\mathbf{A}^{\top} \boldsymbol{\mu} \\
\mathbb{V}\left(\mathbf{A}^{\top} \mathbf{x}\right) & =\mathbf{A}^{\top} \boldsymbol{\Sigma} \mathbf{A} \\
\mathbb{E}\left(\mathbf{x}^{\top} \mathbf{A} \mathbf{x}\right) & =\boldsymbol{\mu}^{\top} \mathbf{A} \boldsymbol{\mu}+\operatorname{tr}(\mathbf{A} \boldsymbol{\Sigma})
\end{aligned}
$$

where $\operatorname{tr}(\mathbf{A})=\sum_{i} A_{i i}$ is the trace of $\mathbf{A}$

- Some useful facts about traces:

$$
\begin{aligned}
\operatorname{tr}(\mathbf{A B}) & =\operatorname{tr}(\mathbf{B} \mathbf{A}) \\
\operatorname{tr}(\mathbf{A}+\mathbf{B}) & =\operatorname{tr}(\mathbf{A})+\operatorname{tr}(\mathbf{B}) \\
\operatorname{tr}(c \mathbf{A}) & =c \operatorname{tr}(\mathbf{A}) \\
\operatorname{tr}(\mathbf{A}) & =\operatorname{rank}(\mathbf{A}) \quad \text { if } \mathbf{A} \mathbf{A}=\mathbf{A}
\end{aligned}
$$

## Eigendecompositions

- Finally, we'll also take a moment to introduce some facts about eigenvalues
- The most important thing about eigenvalues is that they allow us to "diagonalize" a matrix: if $\mathbf{A}$ is a symmetric $d \times d$ matrix, then it can be factored into:

$$
\mathbf{A}=\mathbf{Q} \mathbf{\Lambda} \mathbf{Q}^{\top}
$$

where $\boldsymbol{\Lambda}$ is a diagonal matrix containing the eigenvalues $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{d}$ of $\mathbf{A}$ and the columns of $\mathbf{Q}$ are its eigenvectors

- Furthermore, eigenvectors are orthonormal, so we have $\mathbf{Q}^{\top} \mathbf{Q}=\mathbf{Q Q}^{\top}=\mathbf{I}$


## Eigenvalues and "size"

- This is very helpful from a conceptual standpoint, as it allows us to separate the "size" of a matrix ( $\boldsymbol{\Lambda}$ ) from its "direction(s)" (Q)
- For example, we have already seen that one measure of the size of a matrix is based on $\lambda_{\text {max }}$ (for a symmetric matrix, its spectral norm is its largest eigenvalue)
- In addition, the trace and determinant, two other ways of quantifying the "size" of a matrix, are simple functions of the eigenvalues:
- $\operatorname{tr}(\mathbf{A})=\sum_{i} \lambda_{i}$
- $|\mathbf{A}|=\prod_{i} \lambda_{i}$


## Eigenvalues and inverses

- Once one has obtained the eigendecomposition of $\mathbf{A}$, calculating its inverse is straightforward
- If $\mathbf{A}$ is not singular, then $\mathbf{A}^{-1}=\mathbf{Q} \boldsymbol{\Lambda}^{-1} \mathbf{Q}^{\top}$; note that since $\boldsymbol{\Lambda}$ is diagonal, its inverse is trivial to calculate
- Even if $\mathbf{A}$ is singular, we can obtain something called a "generalized inverse": $\mathbf{A}^{-}=\mathbf{Q} \boldsymbol{\Lambda}^{-} \mathbf{Q}^{\top}$, where $\left(\boldsymbol{\Lambda}^{-}\right)_{i i}=\lambda_{i}^{-1}$ if $\lambda_{i} \neq 0$ and $\left(\boldsymbol{\Lambda}^{-}\right)_{i i}=0$ otherwise
- Many other important properties of matrices can be deduced entirely from their eigenvalues:
- A is positive definite if and only if $\lambda_{i}>0$ for all $i$
- $\mathbf{A}$ is positive semidefinite if and only if $\lambda_{i} \geq 0$ for all $i$
- If $\mathbf{A}$ has rank $r$, then $\mathbf{A}$ has $r$ nonzero eigenvalues and the remaining $d-r$ eigenvalues are zero


## Extreme values

- Lastly, there is a connection between a matrix's eigenvalues and the extreme values of its quadratic form
- Let the eigenvalues $\lambda_{1}, \ldots, \lambda_{d}$ of $\mathbf{A}$ be ordered from largest to smallest. Over the set of all vectors $\mathbf{x}$ such that $\|\mathbf{x}\|_{2}=1$,

$$
\max \mathbf{x}^{\top} \mathbf{A} \mathbf{x}=\lambda_{1}
$$

and

$$
\min \mathbf{x}^{\top} \mathbf{A} \mathbf{x}=\lambda_{d}
$$

