Abstract

This report contains lecture notes for UC Berkeley’s introductory class on Machine Learning. It covers many methods for classification and regression, and several methods for clustering and dimensionality reduction. It is concise because nothing is included that cannot be written or spoken in a single semester’s lectures (with whiteboard lectures and almost no slides!) and because the choice of topics is limited to a small selection of particularly useful, popular algorithms.
Keywords: machine learning, classification, regression, density estimation, dimensionality reduction, clustering, perceptrons, support vector machines (SVMs), Gaussian discriminant analysis, linear discriminant analysis (LDA), quadratic discriminant analysis (QDA), logistic regression, decision trees, random forests, ensemble learning, bagging, boosting, AdaBoost, neural networks, convolutional neural networks (CNNs, ConvNets), nearest neighbor search, least-squares linear regression, logistic regression, polynomial regression, ridge regression, Lasso, bias-variance tradeoff, maximum likelihood estimation (MLE), principal components analysis (PCA), singular value decomposition (SVD), random projection, latent factor analysis, latent semantic indexing, $k$-means clustering, hierarchical clustering, spectral graph clustering, the kernel trick, learning theory
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About this Report

This report compiles my lectures notes for UC Berkeley’s class CS 189/289A, *Machine Learning*, which is both an undergraduate and introductory graduate course. I hope it will serve as a fast introduction to the subject for readers who are already comfortable with vector calculus, linear algebra, probability, and statistics. Please consult my CS 189/289A web page¹ as an addendum to this report; it includes an extended description of each lecture and additional web links and reading assignments related to the lectures. Consider this report and the web page to be living documents; both will be refined a bit every time I teach the class.

The term “lecture notes” has shifted to include long textbook-style treatments written by professors as supplements to their classes. Not so here. This report compiles the actual notes that I lecture from. I call it *Concise Machine Learning* because I include almost nothing that I do not have time to write or speak during one fourteen-week semester of twice-weekly 80-minute lectures. (After holidays and the midterm exam, that amounts to 25 lectures.) Words that appear [in brackets] are spoken; everything else is written on the “whiteboard”—in my class, a tablet computer. My whiteboard software permits me to incorporate (and write on) figures, included here. However, I am largely anti-Powerpoint and I resort to prepared slides for just three or four brief segments during the semester.

These notes might be ideal for mathematically sophisticated readers who want to learn the basics of machine learning as quickly as possible. But they’re not ideal for everybody. The time limitation necessitates that many details are omitted. I believe that the most mathematically well-prepared readers will be able to fill in those details themselves. But many readers, including most students who take the class, will need additional readings or discussion sections for greater detail. My class web page lists additional readings for most of the lectures, many of them from two textbooks that have been kindly made available for free on the web: *An Introduction to Statistical Learning with Applications in R,*² second edition, by Gareth James, Daniela Witten, Trevor Hastie, and Robert Tibshirani, Springer, New York, 2021, ISBN # 978-1-0716-1417-4; and *The Elements of Statistical Learning: Data Mining, Inference, and Prediction,*³ second edition, by Trevor Hastie, Robert Tibshirani, and Jerome Friedman, Springer, New York, 2008. Wikipedia also has excellent introductions to many machine learning algorithms. Readers wanting the verbose kind of “lecture notes” should consider the fine ones written by Stanford University’s Andrew Ng.⁴ I have no interest in duplicating these efforts; instead, I’m aiming for the neglected niche of “shortest introduction.” (And perhaps also “best stolen illustrations.”)

The other thing that makes this report concise is the choice of topics. CS 189/289A was introduced at UC Berkeley in the spring of 2013 by Prof. Jitendra Malik, and most of his topic choices remain intact here. Jitendra told me that he only taught a machine learning algorithm if he or his collaborators had used it successfully for some application. He said, “the machine learning course is too important to leave to the machine learning experts”—that is, users of machine learning algorithms often have a more clear-eyed view of their usefulness than inventors of machine learning algorithms.

I thank Peter Bartlett, Alyosha Efros, Isabelle Guyon, and Jitendra Malik—the previous teachers of CS 189/289A—for their lectures and lecture notes, from which I learned the topic myself. While I’ve given the lectures my own twist and rearranged the material a lot, I am ultimately making incremental improvements (and perhaps incremental worsenings) to a structure they handed down to me. I also thank Carlos Flores for sharing screenshots from my lectures.

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¹https://people.eecs.berkeley.edu/~jrs/189/
²https://www.statlearning.com
³https://hastie.su.domains/ElemStatLearn/
⁴http://cs229.stanford.edu/notes2020spring/
1 Introduction

CS 189 / 289A [Spring 2024]
Machine Learning
Jonathan Shewchuk

https://people.eecs.berkeley.edu/~jrs/189/

Homework 1 due next Wednesday.

Questions: Please use Ed Discussion, not email. [Ed Discussion has an option for private questions, but please use public for most questions so other people can benefit.]
For personal matters only, jrs@berkeley.edu

Discussion sections (Tue & Wed):
Attend any section. [We’ll put up a list on Ed Discussion.]
[We might have a few advanced sections, including research discussion or exam problem preparation.]
Sections start Tuesday. [Next week.]

[Enrollment: 703 students max. 452 waitlisted. Expecting many drops. EECS grads have highest priority; CD/DS undergrads second; non-EECS grads third; a few concurrent enrollment students will be admitted.]

[Textbooks: Available free online. Linked from class web page.]
Prerequisites

Vector calculus: Math 53 [or another vector calculus course]
Linear algebra: Math 54, Math 110, or EE 16A+16B [or another linear algebra course]
Probability: CS 70, EECS 126, or Stat 134 [or another probability course]
Plentiful programming experience [TAs have no obligation to look at your code.]
NOT CS 188

Grading: 189

40% 7 Homeworks. Late policy: 5 slip days total
20% Midterm: TBA, 6:30–8:30 PM
40% Final Exam: Friday, May 10, 3–6 PM

Grading: 289A

40% HW
20% Midterm
20% Final
20% Project

Cheating

– Discussion of HW problems is encouraged. Showing other students small amounts of code is okay.
– All homeworks, including programming, must be written individually.
– We will actively check for plagiarism.
– Typical penalty is a large NEGATIVE score, but I reserve right to give an instant F for even one violation, and will always give an F for two.

[Last time I taught CS 61B, we had to punish roughly 100 people for cheating. It was very painful. Please don’t put me through that again.]

CORE MATERIAL

– Finding patterns in data; using them to make predictions.
– Models and statistics help us understand patterns.
– Optimization algorithms “learn” the patterns.

[The most important part of this is the data. Data drives everything else.
You cannot learn much if you don’t have enough data.
You cannot learn much if your data has bad quality.
But it’s amazing what you can do if you have lots of good data.
Machine learning has changed a lot in the last two decades because the internet has made truly vast quantities of data available. For instance, with a little patience you can download tens of millions of photographs. Then you can build a 3D model of Paris.
Some techniques that had fallen out of favor, like neural networks, have come back big in recent years because researchers found that they work so much better when you have vast quantities of data.]
CLASSIFICATION

- Collect training points with class labels: reliable debtors & defaulted debtors
- Evaluate new applicants—predict their class

![Credit Card Data](creditcardscrop.pdf)

**Introduction**

[The problem of classification. We are given data points, each belonging to one of two classes: orange crosses represent people who defaulted on their credit cards, and blue circles represent those who didn’t. Then we are given additional points whose class is unknown, and we are asked to predict what class each new point is in. Given the credit card balance and annual income of new applicants, predict whether they will default on their debt.]

![Decision Boundaries](classify.pdf)

[Draw this figure by hand.]

[Draw 2 colors of dots, almost but not quite linearly separable.]

[“How do we classify a new point?” Draw a point in a third color.]

[One possibility: look at its nearest neighbor.]

[Another possibility: draw a linear decision boundary; label it.]

[Those are two different models for the nature of this data.]
[We’ll learn some ways to compute linear decision boundaries in the next several lectures. But for now, let’s compare these two methods.]

[Two examples of classifiers for the same data: a nearest neighbor classifier (left) and a linear classifier (right). The decision boundaries are in black.]

[At left we have a nearest neighbor classifier, which classifies a new point by finding the nearest point in the training data, and assigning it the same class. At right we have a linear classifier, which guesses that everything above the line is brown, and everything below the line is blue. At right, the linear decision boundary—the black line—is explicitly computed by the classifier. At left, the decision boundary is not computed; the classifier just takes a new point and computes the distances to all the training points.]

[The neighbor classifier at left has a big advantage: it classifies all the training data correctly, whereas the linear classifier does not. But the linear classifier has an advantage too. Somebody please tell me what.]

[A 1-nearest neighbor classifier and a 15-nearest neighbor classifier.]

[The 15-nearest neighbor classifier classifies a new point by looking at its 15 nearest neighbors and letting them vote for the correct class.]

[The left figure is an example of what’s called overfitting. In the left figure, observe how intricate the decision boundary is that separates the positive examples from the negative examples. It’s a bit too intricate to reflect reality. In the right figure, the decision boundary is smoother. Intuitively, that smoothness is probably more likely to correspond to reality.]
Classifying Digits

In the MNIST digit recognition problem, we are given handwritten digits, and we are asked to learn to distinguish them. See Homework 1.

Express these images as vectors

\[
\begin{bmatrix}
3 & 3 & 3 & 3 \\
0 & 0 & 2 & 3 \\
0 & 0 & 1 & 3 \\
3 & 3 & 3 & 3 \\
\end{bmatrix}
\rightarrow
\begin{bmatrix}
3 \\
3 \\
3 \\
3 \\
\end{bmatrix}
\]

Images are points in 16-dimensional space. Linear decision boundary is a hyperplane.

TRAIN, VALIDATE, TEST

How we classify:

- We are given labeled data—sample points with class labels.
- Hold back a subset of the labeled points, called the validation set. Maybe 20%. The other 80% is the training set. [Warning: the term training data is not used consistently. Often “training data” refers to all the labeled data. You have to judge from context.]
- Train one or more classifiers; they learn to distinguish 7 from not 7. Use training set to learn model weights. Do NOT use validation set to train!!!
- Usually, train multiple learning algorithms, or one algorithm with multiple hyperparameter settings, or both [using the same training set for each].
- Validate the trained classifiers on the validation set. Choose classifier/hyperparameters with lowest validation error. Called validation. [When we do validation, we are not learning any more. We are checking what classes our trained classifiers assign to our validation set, and counting how often they’re right. We use this to judge our models—not how well they remember the training set labels.]
- Optional: Test the best classifier on a test set of NEW data. Final evaluation. Typically you do NOT have the labels. [But somebody else might have them, and assign you a score!]
3 kinds of error:

- **Training error**: fraction of training set not classified correctly. [This is zero with the 1-nearest neighbor classifier, but nonzero with the 15-nearest neighbor and linear classifiers. But that doesn’t mean the 1-nearest neighbor classifier is always better. Remember that you cannot include the validation data in this calculation, even if somebody calls it “training data.”]

- **Validation error**: fraction of validation set misclassified. Use this to choose classifier/hyperparameters. [You didn’t use the validation set to train, so even the 1-nearest neighbor classifier can classify these points wrong. Validation error is almost always higher than training error.]

- **Test error**: fraction of test set misclassified. Used to evaluate you.

Most ML algorithms have a few hyperparameters that control over/underfitting, e.g. $k$ in $k$-nearest neighbors.

overfitlabeled.pdf (modified from ESL, Figure 2.4)

- **overfitting**: when the validation/test error deteriorates because the classifier becomes too sensitive to outliers or other spurious patterns.

- **underfitting**: when the validation/test error deteriorates because the classifier is not flexible enough to fit patterns.

- **outliers**: points with atypical labels (e.g., rich borrower who defaulted anyway). Increase risk of overfitting.

[In machine learning, the goal is to create a classifier that generalizes to new examples we haven’t seen yet. Overfitting and underfitting are both counterproductive to that goal. So we’re always seeking a compromise: we want decision boundaries that make fine distinctions without being downright superstitious.]

Kaggle.com:

- Runs ML competitions, including our HWs
- We use 2 data sets:
  - public set labels available during competition
  - private test set labels known only to Kaggle

[If your public results are a lot better than your private results, we will know that you overfitted.]
2 Linear Classifiers and Perceptrons

CLASSIFIERS

You are given sample of \( n \) observations, each with \( d \) features [aka predictors]. Some observations belong to class \( C \); some do not.

Example: Observations are bank loans
Features are income & age \((d = 2)\)
Some are in class “defaulted,” some are not

Goal: Predict whether future borrowers will default, based on their income & age.

Represent each observation as a point in \( d \)-dimensional space, called a sample point / a feature vector / independent variables.

We draw these lines/curves separating \( C \)’s from \( X \)’s. Then we use these curves to predict which future borrowers will default. In the last example, though, we’re probably overfitting, which could hurt our predictions.

decision boundary: the boundary chosen by our classifier to separate items in the class from those not.

overfitting: When sinuous decision boundary fits sample points so well that it doesn’t classify future points well.

[A reminder that underlined phrases are definitions, worth memorizing.]

Some (not all) classifiers work by computing a

decision function: A function \( f(x) \) that maps a point \( x \) to a scalar such that

\[
\begin{align*}
f(x) > 0 & \quad \text{if } x \in \text{class } C; \\
f(x) \leq 0 & \quad \text{if } x \notin \text{class } C.
\end{align*}
\]

Aka predictor function or discriminant function.

For these classifiers, the decision boundary is \( \{x \in \mathbb{R}^d : f(x) = 0\} \)
[That is, the set of all points where the decision function is zero.]
Usually, this set is a \((d - 1)\)-dimensional surface in \( \mathbb{R}^d \).

\( \{x : f(x) = 0\} \) is also called an isosurface of \( f \) for the isovalue \( 0 \).

\( f \) has other isosurfaces for other isovalue, e.g., \( \{x : f(x) = 1\} \).
[3D plot and isocontour plot of the cone] $f(x, y) = \sqrt{x^2 + y^2} - 3$.

[Imagine a decision function in $\mathbb{R}^d$, and imagine its $(d - 1)$-dimensional isosurfaces.]

[One of these spheres could be the decision boundary.]

**linear classifier**: The decision boundary is a line/plane. Usually uses a linear decision function. [Sometimes no decision fn.]
Math Review

[I will write vectors in matrix notation.]

Vectors: \[ x = \begin{bmatrix} x_1 \\ x_2 \\ x_3 \\ x_4 \\ x_5 \end{bmatrix} = [x_1 \ x_2 \ x_3 \ x_4 \ x_5]^T \]

Think of \( x \) as a point in 5-dimensional space.

Conventions (often, but not always):
- uppercase roman = matrix, random variable, set \( X \)
- lowercase roman = vector \( x \)
- Greek = real scalar \( \alpha \)
- Some integers:
  - \( n \) = # of sample points
  - \( d \) = # of features (per point)
  - = dimension of sample points
  - \( i \ j \ k \) = indices

Function (often scalar)
- \( f(\ ), s(\ ), \ldots \)

Euclidean inner product (aka dot product):
- \( x \cdot y = x_1y_1 + x_2y_2 + \ldots + x_dy_d \)
- Also written \( x^\top y \)

Clearly, \( f(x) = w \cdot x + \alpha \) is a linear function in \( x \).

Euclidean norm:
- \( ||x|| = \sqrt{x \cdot x} = \sqrt{x_1^2 + x_2^2 + \ldots + x_d^2} \)
- \( ||x|| \) is the length (aka Euclidean length) of a vector \( x \).
- Given a vector \( x \neq 0 \), \( \frac{x}{||x||} \) is a unit vector (length 1).
- “Normalize a vector \( x \)”: replace \( x \) with \( \frac{x}{||x||} \).

Use dot products to compute angles:

\[
\cos \theta = \frac{x \cdot y}{||x|| ||y||} = \frac{x}{||x||} \cdot \frac{y}{||y||} \quad \text{(length 1 length 1)}
\]

\[
\begin{align*}
&\text{acute} & x \cdot y > 0 \\
&\text{right} & x \cdot y = 0 \\
&\text{obtuse} & x \cdot y < 0
\end{align*}
\]

Given a linear decision function \( f(x) = w \cdot x + \alpha \), the decision boundary is

\[ H = \{ x : w \cdot x = -\alpha \} \]

The set \( H \) is called a hyperplane. (A line in 2D, a plane in 3D.)

[A hyperplane is what you get when you generalize the idea of a plane to higher dimensions. The three most important things to understand about a hyperplane is (1) it has dimension \( d - 1 \) and it cuts the \( d \)-dimensional space into two halves; (2) it’s flat; and (3) it’s infinite.]
Theorem: Let $x, y$ be 2 points that lie on $H$. Then $w \cdot (y - x) = 0$.

Proof: $w \cdot (y - x) = -\alpha - (-\alpha) = 0$. Therefore, $w$ is orthogonal to any line segment that lies on $H$. $w$ is called the normal vector of $H$, because (as the theorem shows) $w$ is normal (perpendicular) to $H$. [I.e., $w$ is perpendicular to every line on $H$.]

If $w$ is a unit vector, then $f(x) = w \cdot x + \alpha$ is the signed distance from $x$ to $H$. [See Discussion 1.] I.e., positive on $w$’s side of $H$; negative on other side.

Moreover, the distance from $H$ to the origin is $\alpha$. [How do we know that?] Hence $\alpha = 0$ if and only if $H$ passes through origin.

[w does not have to be a unit vector for the classifier to work.
If $w$ is not a unit vector, $w \cdot x + \alpha$ is the signed distance times some real.
If you want to fix that, you can rescale the equation by computing $\|w\|$ and dividing both $w$ and $\alpha$ by $\|w\|$.

The coefficients in $w$, plus $\alpha$, are called weights (or parameters or regression coefficients).

[That’s why we call the vector $w$; “$w$” stands for “weights.”]

The training points are linearly separable if there exists a hyperplane that correctly classifies all the training points.

[At the beginning of this lecture, I showed you one plot that’s linearly separable and two that are not.]

[We will investigate some linear classifiers that only work for linearly separable data and some that do a decent job with non-separable data. Obviously, if your data are not linearly separable, a linear classifier cannot do a perfect job. But we’re still happy if we can find a classifier that usually predicts correctly.]

**A Simple Classifier**

**Centroid method**: compute mean $\mu_C$ of all training points in class $C$ and mean $\mu_X$ of all points NOT in $C$.

We use the decision function

$$f(x) = \frac{(\mu_C - \mu_X) \cdot x - (\mu_C - \mu_X) \cdot \mu_C + \mu_X}{2}$$

so the decision boundary is the hyperplane that bisects line segment w/endpoint $\mu_C, \mu_X$. 

In this example, there’s clearly a linear classifier that classifies every training point correctly, and the centroid method isn’t it.

Note that this is hardly the worst example I could have given.

If you’re in the mood for an easy puzzle, pull out a sheet of paper and think of an example, with lots of training points, where the centroid method misclassifies every training point but two.

Nevertheless, there are circumstances where this method works well, like when all your positive examples come from one Gaussian distribution, and all your negative examples come from another.

We can sometimes improve this classifier by adjusting the scalar term $\alpha$ to minimize the number of misclassified points. Then the hyperplane has the same normal vector, but a different position.

**Perceptron Algorithm (Frank Rosenblatt, 1957)**

Slow, but correct for linearly separable points.

Uses a numerical optimization algorithm, namely, gradient descent.

Poll:
- How many of you know what gradient descent is?
- How many of you know what a linear program is?
- How many of you know what the simplex algorithm for linear programming is?
- How many of you know what a quadratic program is?

We’re going to learn what most of these things are. As machine learning people, we will be heavy users of optimization methods. Unfortunately, I won’t have time to teach you algorithms for many optimization problems, but we’ll learn a few. To learn more, take EECS 127.

Consider $n$ sample points $X_1, X_2, ..., X_n$.

The reason I’m using capital $X$ here is because we typically store these vectors as rows of a matrix $X$. So the subscript picks out a row of $X$, representing a specific sample point.

For each sample point, the label $y_i = \begin{cases} 1 \text{ if } X_i \in \text{class } C, \\ -1 \text{ if } X_i \not\in C. \end{cases}$

For simplicity, consider only decision boundaries that pass through the origin. (We’ll fix this later.)
Goal: find weights $w$ such that
\[ X_i \cdot w \geq 0 \quad \text{if } y_i = 1, \quad \text{and} \]
\[ X_i \cdot w \leq 0 \quad \text{if } y_i = -1. \quad \text{[remember, } X_i \cdot w \text{ is the signed distance]} \]
Equivalently: $y_i X_i \cdot w \geq 0$. ← inequality called a constraint.

Idea: We define a risk function $R$ that is positive if some constraints are violated. Then we use optimization to choose $w$ that minimizes $R$. [That’s how we train a perceptron classifier.]

Define the loss function
\[
L(z, y_i) = \begin{cases} 
0 & \text{if } y_i z \geq 0, \\
-y_i z & \text{otherwise.}
\end{cases}
\]
[Here, $z$ is the classifier’s prediction, and $y_i$ is the correct answer.]

If $z$ has the same sign as $y_i$, the loss function is zero (happiness).
But if $z$ has the wrong sign, the loss function is positive.

[For each training point, you want to get the loss function down to zero, or as close to zero as possible. It’s called the “loss function” because the bigger it is, the bigger a loser your classifier is.]

Define risk function (aka objective function or cost function)
\[
R(w) = \frac{1}{n} \sum_{i=1}^{n} L(X_i \cdot w, y_i)
\]
\[
= \frac{1}{n} \sum_{i \in V} -y_i X_i \cdot w \quad \text{where } V \text{ is the set of indices } i \text{ for which } y_i X_i \cdot w < 0.
\]

If $w$ classifies all $X_1, \ldots, X_n$ correctly, then $R(w) = 0$.
Otherwise, $R(w)$ is positive, and we want to find a better $w$.

Goal: Solve this optimization problem:
\[
\text{Find } w \text{ that minimizes } R(w).
\]

[Plot of risk $R(w)$. Every point in the dark green flat spot is a minimum. We’ll look at this more next lecture.]
3 Perceptron Learning; Maximum Margin Classifiers

Perceptron Algorithm (cont’d)

Recall:
– linear decision fn \( f(x) = w \cdot x \) (for simplicity, no \( \alpha \))
– decision boundary \( \{ x : f(x) = 0 \} \) (a hyperplane through the origin)
– sample points \( X_1, X_2, \ldots, X_n \in \mathbb{R}^d \), class labels \( y_1, \ldots, y_n = \pm 1 \)
– goal: find weights \( w \) such that \( y_i X_i \cdot w \geq 0 \)
– goal, revised: find \( w \) that minimizes \( R(w) = \sum_{i \in V} -y_i X_i \cdot w \) [risk function]

where \( V = \{ i : y_i X_i \cdot w < 0 \} \).

[Our original problem was to find a separating hyperplane in one space, which I’ll call \( x \)-space. But we’ve transformed this into a problem of finding an optimal point in a different space, which I’ll call \( w \)-space. It’s important to understand transformations like this, where a geometric structure in one space becomes a point in another space.]

Objects in \( x \)-space transform to objects in \( w \)-space:

\[
\begin{array}{|c|c|}
\hline
\text{\( x \)-space} & \text{\( w \)-space} \\
\hline
\text{hyperplane: } \{ z : w \cdot z = 0 \} & \text{point: } \boldsymbol{w} \\
\text{point: } \boldsymbol{x} & \text{hyperplane: } \{ z : x \cdot z = 0 \} \\
\hline
\end{array}
\]

Point \( x \) lies on hyperplane \( \{ z : w \cdot z = 0 \} \) \( \iff \) \( w \cdot x = 0 \) \( \iff \) point \( w \) lies on hyperplane \( \{ z : x \cdot z = 0 \} \) in \( w \)-space.

[So a hyperplane transforms to a point that represents its normal vector. And a sample point transforms to the hyperplane whose normal vector is the sample point.]

[In this algorithm, the transformations happen to be symmetric: a hyperplane in \( x \)-space transforms to a point in \( w \)-space the same way that a hyperplane in \( w \)-space transforms to a point in \( x \)-space. That won’t always be true for the decision boundaries we use this semester.]

If we want to enforce inequality \( x \cdot w \geq 0 \), that means
– in \( x \)-space, \( x \) should be on the same side of \( \{ z : w \cdot z = 0 \} \) as \( w \)
– in \( w \)-space, \( w \) is on the same side of \( \{ z : x \cdot z = 0 \} \) as \( x \)

[Draw this by hand. xwspace.pdf]
[Observe that the \( x \)-space sample points are the normal vectors for the \( w \)-space lines. We can choose \( w \) to be anywhere in the shaded region.]

[For a sample point \( x \) in class \( C \), \( w \) and \( x \) must be on the same side of the hyperplane that \( x \) transforms into. For a point \( x \) not in class \( C \) (marked by an \( X \)), \( w \) and \( x \) must be on opposite sides of the hyperplane that \( x \) transforms into. These rules determine the shaded region above, in which \( w \) must lie.]

[Again, what have we accomplished? We have switched from the problem of finding a hyperplane in \( x \)-space to the problem of finding a point in \( w \)-space. That’s a better fit to how we think about optimization algorithms.]
[Let’s take a look at the risk function these three sample points create.]

\[ \text{riskplot.pdf, riskiso.pdf} \] [Plot & isocontours of risk \( R(w) \). Note how \( R \)'s creases match the lines in the \( w \)-space drawn above.]

[In this plot, we can choose \( w \) to be any point in the bottom pizza slice; all those points minimize \( R \).]

[We have an optimization problem; we need an optimization algorithm to solve it.]

An optimization algorithm: gradient descent on \( R \).

[Draw the typical steps of gradient descent on the plot of \( R \).]

Given a starting point \( w \), find gradient of \( R \) with respect to \( w \); this is the direction of steepest ascent. Take a step in the opposite direction. Recall [from your vector calculus class]

\[
\nabla R(w) = \begin{bmatrix} \frac{\partial R}{\partial w_1} \\ \frac{\partial R}{\partial w_2} \\ \vdots \\ \frac{\partial R}{\partial w_d} \end{bmatrix} \quad \text{and} \quad \nabla w(z \cdot w) = \begin{bmatrix} z_1 \\ z_2 \\ \vdots \\ z_d \end{bmatrix} = z
\]

\[
\nabla R(w) = \sum_{i \in V} \nabla - y_i X_i \cdot w = - \sum_{i \in V} y_i X_i
\]

At any point \( w \), we walk downhill in direction of steepest descent, \(-\nabla R(w)\).

\[
w \leftarrow \text{arbitrary nonzero starting point (good choice is any } y_i X_i \text{)}
\]

while \( R(w) > 0 \)

\[
V \leftarrow \text{set of indices } i \text{ for which } y_i X_i \cdot w < 0
\]

\[
w \leftarrow w + \epsilon \sum_{i \in V} y_i X_i
\]

return \( w \)

\( \epsilon > 0 \) is the step size aka learning rate, chosen empirically. [Best choice depends on input problem!]

Problem: Slow! Each step takes \( O(nd) \) time. [Can we improve this?]
Optimization algorithm 2: stochastic gradient descent

Idea: each step, pick one misclassified $X_i$; do gradient descent on loss fn $L(X_i \cdot w, y_i)$.

Called the perceptron algorithm. Each step takes $O(d)$ time. [Not counting the time to search for a misclassified $X_i$.]

\[
\text{while some } y_i X_i \cdot w < 0 \\
\quad w \leftarrow w + \epsilon y_i X_i \\
\text{return } w
\]

[Stochastic gradient descent is quite popular and we’ll see it several times more this semester, especially for neural networks. However, stochastic gradient descent does not work for every problem that gradient descent works for. The perceptron risk function happens to have special properties that guarantee that stochastic gradient descent will always succeed.]

What if separating hyperplane doesn’t pass through origin? Add a fictitious dimension. Decision fn is

\[
f(x) = w \cdot x + \alpha = [w_1 \ w_2 \ \alpha] \cdot \begin{bmatrix} x_1 \\ x_2 \\ 1 \end{bmatrix}
\]

Now we have sample points in $\mathbb{R}^{d+1}$, all lying on hyperplane $x_{d+1} = 1$.

Run perceptron algorithm in $(d + 1)$-dimensional space. [We are simulating a general hyperplane in $d$ dimensions by using a hyperplane through the origin in $d + 1$ dimensions.]

[The perceptron algorithm was invented in 1957 by Frank Rosenblatt at the Cornell Aeronautical Laboratory. It was originally designed not to be a program, but to be implemented in hardware for image recognition on a 20 × 20 pixel image. Rosenblatt built a Mark I Perceptron Machine that ran the algorithm, complete with electric motors to do weight updates.]
[Then he held a press conference where he predicted that perceptrons would be “the embryo of an electronic computer that [the Navy] expects will be able to walk, talk, see, write, reproduce itself and be conscious of its existence.” We’re still waiting on that.]

[One interesting aspect of the perceptron algorithm is that it’s an “online algorithm,” which means that if new training points come in while the algorithm is already running, you can just throw them into the mix and keep looping.]

[Perceptron Convergence Theorem: If data is linearly separable, perceptron algorithm will find a linear classifier that classifies all data correctly in at most $O(r^2/\gamma^2)$ iterations, where $r = \max \|X_i\|$ is “radius of data” and $\gamma$ is the “maximum margin.”]

[I’ll define “maximum margin” shortly.]

[We’re not going to prove this, because perceptrons are obsolete.]

[Although the step size/learning rate $\epsilon$ doesn’t appear in that big-O expression, it does have an effect on the running time, but the effect is hard to characterize. The algorithm gets slower if $\epsilon$ is too small because it has to take lots of steps to get down the hill. But it also gets slower if $\epsilon$ is too big for a different reason: it jumps right over the region with zero risk and oscillates back and forth for a long time.]

[Although stochastic gradient descent is faster for this problem than gradient descent, the perceptron algorithm is still slow. There’s no reliable way to choose a good step size $\epsilon$. Fortunately, optimization algorithms have improved a lot since 1957. You can get rid of the step size by using a decent modern “line search” algorithm. Better yet, you can find a better decision boundary much more quickly by quadratic programming, which is what we’ll talk about next.]

**MAXIMUM MARGIN CLASSIFIERS**

The margin of a linear classifier is the distance from the decision boundary to the nearest training point. What if we make the margin as wide as possible?

![Diagram of maximum margin classifiers](maxmargin.pdf)

We enforce the constraints

$$y_i (w \cdot X_i + \alpha) \geq 1 \quad \text{for } i \in [1, n]$$

[Notice that the right-hand side is a 1, rather than a 0 as it was for the perceptron algorithm. It’s not obvious, but this a better way to formulate the problem, partly because it makes it impossible for the weight vector $w$ to get set to zero.]
Recall: if $\|w\| = 1$, signed distance from hyperplane to $X_i$ is $w \cdot X_i + \alpha$. [We’ve normalized the expression to get a unit weight vector.]

Otherwise, it’s $\frac{w}{\|w\|} \cdot X_i + \frac{\alpha}{\|w\|}$. [We’ve normalized the expression to get a unit weight vector.]

Hence the margin is $\min_i \frac{1}{\|w\|} |w \cdot X_i + \alpha| \geq \frac{1}{\|w\|}$. [We get the inequality by substituting the constraints.]

To maximize the margin, minimize $\|w\|$. Optimization problem:

| Find $w$ and $\alpha$ that minimize $\|w\|^2$ |
| subject to $y_i(X_i \cdot w + \alpha) \geq 1$ for all $i \in [1, n]$ |

Called a quadratic program in $d + 1$ dimensions and $n$ constraints.

It has one unique solution! [If the points are linearly separable; otherwise, it has no solution.]

[A reason we use $\|w\|^2$ as an objective function, instead of $\|w\|$, is that the length function $\|w\|$ is not smooth at $w = 0$, whereas $\|w\|^2$ is smooth everywhere. This makes optimization easier.]

The solution gives us a maximum margin classifier, aka a hard-margin support vector machine (SVM).

[Technically, this isn’t really a support vector machine yet; it doesn’t fully deserve that name until we add features and kernels, which we’ll study in later lectures.]

At the optimal solution, the margin is exactly $\frac{1}{\|w\|}$. [Because at least one constraint holds with equality.]

There is a slab of width $\frac{2}{\|w\|}$ containing no sample points [with the hyperplane running along its middle]. [Let’s see what these constraints look like in weight space.]

This is an example of what the linear constraints look like in the 3D weight space $(w_1, w_2, \alpha)$ for the SVM we’ve been studying with three training points. The SVM is looking for the point nearest the $\alpha$-axis that lies above the blue plane (representing an in-class training point) but below the red and pink planes (representing out-of-class training points). In this example, that optimal point lies where the three planes intersect. At right we see a 2D cross section $w_1 = 1/17$ of the 3D space, because the optimal solution lies in this cross section. The constraints say that the solution must lie in the leftmost pizza slice, while being as close to the origin as possible, so the optimal solution is where the three lines meet.

[Like the perceptron algorithm, a hard-margin SVM works only with linearly separable point sets. We’ll fix that in the next lecture.]
4 Soft-Margin Support Vector Machines; Features

SOFT-MARGIN SUPPORT VECTOR MACHINES (SVMs)

Solves 2 problems:
- Hard-margin SVMs fail if data not linearly separable.
- “ ” “ sensitive to outliers.

[Figure 9.5. Left: Two classes of observations are shown in blue and in purple, along with the maximal margin hyperplane. Right: An additional blue observation has been added, leading to a dramatic shift in the maximal margin hyperplane shown as a solid line. The dashed line indicates the maximal margin hyperplane that was obtained in the absence of this additional point.]

• Greater robustness to individual observations, and
• Better classification of most of the training observations.

That is, it could be worthwhile to misclassify a few training observations in order to do a better job in classifying the remaining observations.

The support vector classifier, sometimes called a soft margin classifier, does exactly this. Rather than seeking the largest possible margin so that every observation is not only on the correct side of the hyperplane but also on the correct side of the margin, we instead allow some observations to be on the incorrect side of the margin, or even the incorrect side of the hyperplane. (The margin is soft because it can be violated by some of the training observations.) An example is shown in the left-hand panel of Figure 9.6. Most of the observations are on the correct side of the margin. However, a small subset of the observations are on the wrong side of the margin.

An observation can be not only on the wrong side of the margin, but also on the wrong side of the hyperplane. In fact, when there is no separating hyperplane, such a situation is inevitable. Observations on the wrong side of the hyperplane correspond to training observations that are misclassified by the support vector classifier. The right-hand panel of Figure 9.6 illustrates such a scenario.

9.2.2 Details of the Support Vector Classifier

The support vector classifier classifies a test observation depending on which side of a hyperplane it lies. The hyperplane is chosen to correctly classify a test observation depending on which side of a hyperplane it lies. The hyperplane is chosen to correctly

Idea: Allow some points to violate the margin, with slack variables.

Modified constraint for point $i$:

$$y_i (X_i \cdot w + \alpha) \geq 1 - \xi_i$$

[Observe that the only difference between these constraints and the hard-margin constraints we saw last lecture is the extra slack term $\xi_i$.]

[We also impose new constraints, that the slack variables are never negative.]

$$\xi_i \geq 0$$

[This inequality ensures that all sample points that don’t violate the margin are treated the same; they all have $\xi_i = 0$. Point $i$ has nonzero $\xi_i$ if and only if it violates the margin.]

Re-define “margin” to be $1/\|w\|$. [For soft-margin SVMs, the margin is no longer the distance from the decision boundary to the nearest training point; instead, it’s $1/\|w\|$.]
To prevent abuse of slack, we add a **loss term** to objective fn.

**Optimization problem:**

Find $w, \alpha$, and $\xi_i$ that minimize $\|w\|^2 + C \sum_{i=1}^{n} \xi_i$
subject to $y_i(x_i \cdot w + \alpha) \geq 1 - \xi_i$ for all $i \in [1, n]$
$\xi_i \geq 0$ for all $i \in [1, n]$

... a quadratic program in $d + n + 1$ dimensions and $2n$ constraints.

[It's a quadratic program because its objective function is quadratic and its constraints are linear inequalities.]

$C > 0$ is a scalar regularization hyperparameter that trades off:

<table>
<thead>
<tr>
<th></th>
<th>small C</th>
<th>big C</th>
</tr>
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<tbody>
<tr>
<td><strong>desire</strong></td>
<td>maximize margin $1/|w|$</td>
<td>keep most slack variables zero or small</td>
</tr>
<tr>
<td><strong>danger</strong></td>
<td>underfitting (misclassifies much training data)</td>
<td>overfitting (awesome training, awful test)</td>
</tr>
<tr>
<td><strong>outliers</strong></td>
<td>less sensitive</td>
<td>very sensitive</td>
</tr>
<tr>
<td><strong>boundary</strong></td>
<td>more “flat”</td>
<td>more sinuous</td>
</tr>
</tbody>
</table>

[The last row only applies to nonlinear decision boundaries, which we’ll discuss next. Obviously, a linear decision boundary can’t be “sinuous.”]

Use validation to choose $C$.

---

[Examples of how the slab varies with $C$. Smallest $C$ at upper left; largest $C$ at lower right.]
**FEATURES**

Q: How to do nonlinear decision boundaries?

A: Make nonlinear features that lift points into a higher-dimensional space.

High-$d$ linear classifier $\rightarrow$ low-$d$ nonlinear classifier.

Features work with all classifiers—not only linear classifiers like perceptrons and SVMs, but also classifiers that are not linear.

**Example 1: The parabolic lifting map**

\[
\Phi : \mathbb{R}^d \rightarrow \mathbb{R}^{d+1}
\]

\[
\Phi(x) = \begin{bmatrix} x \\ \|x\|^2 \end{bmatrix} \quad \text{left lifts } x \text{ onto paraboloid } x_{d+1} = \|x\|^2
\]

[We’ve added one new feature, $\|x\|^2$. Even though the new feature is just a function of other input features, it gives our linear classifier more power. Now an SVM can have spheres as decision boundaries.]

Find a linear classifier in $\Phi$-space.

It induces a sphere classifier in $x$-space.

Theorem: $\Phi(X_1), \ldots, \Phi(X_n)$ are linearly separable iff $X_1, \ldots, X_n$ are separable by a hypersphere.

(Possibly an $\infty$-radius hypersphere $= \text{hyperplane}$.)

Proof: Consider hypersphere in $\mathbb{R}^d$ w/center $c$ & radius $\rho$. $x$ is inside iff

\[
\|x - c\|^2 < \rho^2
\]

\[
\|x\|^2 - 2c \cdot x + \|c\|^2 < \rho^2
\]

\[
\begin{bmatrix} -2c^T \\ 1 \end{bmatrix} \begin{bmatrix} x \\ \|x\|^2 \end{bmatrix} < \rho^2 - \|c\|^2
\]

Hence points inside sphere $\leftrightarrow$ lifted points underneath hyperplane in $\Phi$-space.

[The implication works in both directions.]

[Hyperspheres include hyperplanes as a special, degenerate case. A hyperplane is essentially a hypersphere with infinite radius. So hypersphere decision boundaries can do everything hyperplane decision boundaries can do, plus a lot more. With the parabolic lifting map, if you pick a hyperplane in $\Phi$-space that is vertical, you get a hyperplane in $x$-space.]
**Example 2: Ellipsoid/hyperboloid/paraboloid decision boundaries**

[Draw 2D examples of ellipse & hyperbola.]

In 3D, these have the formula

\[ Ax_1^2 + Bx_2^2 + Cx_3^2 + Dx_1x_2 + Ex_2x_3 + Fx_3x_1 + Gx_1 + Hx_2 + Ix_3 + \alpha = 0 \]

[Here, the capital letters are scalars, not matrices.]

![quadrics.png (courtesy Rahul Narain)](quadrics.png)

[Quadrics in 3D.] If we add all the quadratic monomials as features, our decision boundaries can be arbitrary ellipsoids, hyperboloids, and paraboloids.

\[
\Phi(x) = [x_1^2 \quad x_2^2 \quad x_3^2 \quad x_1x_2 \quad x_2x_3 \quad x_3x_1 \quad x_1 \quad x_2 \quad x_3]^T
\]

[For perceptron or regression, add a 1 at end. For SVM, the 1 is built-in.]

Decision function is

\[
\begin{bmatrix}
A & B & C & D & E & F & G & H & I
\end{bmatrix} \cdot \Phi(x) + \alpha
\]

[Now, our decision function can be any degree-2 polynomial.]

Isosurface defined by this equation is called a quadric.

A linear decision boundary in \( \Phi \)-space imposes a quadric decision boundary in \( x \)-space.

[The word *quadric* just means an isosurface of a degree-2 polynomial. In the special case of two dimensions, it’s also known as a *conic section*. Our decision boundary can be an arbitrary ellipsoid, hyperboloid, or paraboloid.]

[When \( d \) is large, there are order-\( d^2 \) cross-terms in \( \Phi \)-space! So we are adding a lot of new features. This will impose a serious computational cost on a classifier like a support vector machine. But it might be worth it to find good classifiers for data that aren’t linearly separable.]

\[
\Phi(x) : \mathbb{R}^d \rightarrow \mathbb{R}^{(d^2+3d)/2}
\]

[For perceptron or regression, add 1 for the fictitious dimension.]

[If all these extra features make the classifier overfit or make it too slow, you can leave out the cross-terms and include only quadratic terms like \( x_1^2, x_2^2 \), etc. Then the number of added features is linear in \( d \), not quadratic in \( d \). If you do that, your decision boundaries can be axis-aligned ellipsoids and axis-aligned hyperboloids, but they can’t be rotated in arbitrary ways.]
Example 3: Decision fn is degree-$p$ polynomial

E.g., a cubic in $\mathbb{R}^2$:

$$\Phi(x) = [x_1^3, x_1^2x_2, x_1x_2^2, x_2^3, x_1^2x_2, x_1x_2, x_2^2, x_1, x_2]^T$$

$$\Phi(x) : \mathbb{R}^d \rightarrow \mathbb{R}^{O(d^p)}$$

[Now we’re really blowing up the number of features! If you have, say, 100 features per sample point and you want to use degree-4 decision functions, then each lifted feature vector has a length of roughly 4 million, and your learning algorithm will take approximately forever to run.]

[However, later in the semester we will learn an extremely clever trick that allows us to work with these huge feature vectors very quickly, without ever computing them. It’s called “kernelization” or “the kernel trick.” So even though it appears now that working with degree-4 polynomials is computationally infeasible, it can actually be done quickly.]

Increasing the degree like this accomplishes two things.

- First, the data might become linearly separable when you lift them to a high enough degree, even if the original data are not linearly separable.
- Second, raising the degree can widen the margin, so you might get a more robust decision boundary that generalizes better to test data.

However, if you raise the degree too high, you will overfit the data and then generalization will get worse.]
Features

SVMs with polynomial features of various degrees:
- \( \text{degree} = 1 \)
- \( \text{degree} = 2 \)
- \( \text{degree} = 5 \)

The features we choose are very important! We want them to be rich enough to accurately represent a good classifier. This suggests we should make our set of features as rich as possible:
- linear is a special case of quadratic
- quadratic is a special case of cubic...

Why not include them all? The richer the set of features, the more likely we will encounter overfitting. It's a balancing act: we want our features to be as complex as necessary to represent the classifier, but no more complex.

Features and overfitting

What happens to this picture as sample size grows?

Digit recognition
What features should we use?
- Grey scale level for each pixel
- Orientation histograms
- Polynomials of these

Training vs. test error for degree 1/2/5 decision functions. (Artist’s conception; these aren’t actual calculations, just hand-drawn guesses. Please send me email if you know where to find figures like this with actual data.) In this example, a degree-2 decision gives the smallest test error.

You should search for the ideal degree—not too small, not too big. It’s a balancing act between underfitting and overfitting. The degree is an example of a hyperparameter that can be optimized by validation.

If you’re using both polynomial features and a soft-margin SVM, now you have two hyperparameters: the degree and the regularization hyperparameter \( C \). Generally, the optimal \( C \) will be different for every polynomial degree, so when you change the degree, you should run validation again to find the best \( C \) for that degree.
[So far I’ve talked only about polynomial features. But features can get much more complicated than polynomials, and they can be tailored to fit a specific problem. Let’s consider a type of feature you might use if you wanted to implement, say, a handwriting recognition algorithm.]

**Example 5: Edge detection**

Edge detector: algorithm for approximating grayscale/color gradients in image, e.g.,
- tap filter
- Sobel filter
- oriented Gaussian derivative filter
  [images are discrete, not continuous fields, so approximation of gradients is necessary.]

[See “Image Derivatives” on Wikipedia.]
Collect line orientations in local histograms (each having 12 orientation bins per region); use histograms as features (*instead* of raw pixels).


[If you want to, optionally, use these features in future homeworks and try to win the Kaggle competition, this paper is a good online resource.]

[When they use a linear SVM on the raw pixels, Maji & Malik get an error rate of 15.38% on the test set. When they use a linear SVM on the histogram features, the error rate goes down to 2.64%.]

[Many applications can be improved by designing application-specific features. There’s no limit but your own creativity and ability to discern the structure hidden in your application.]
### 5 Machine Learning Abstractions and Numerical Optimization

**ML ABSTRACTIONS** [some meta comments on machine learning]

When you write a large computer program, you break it down into subroutines and modules. Many of you know from experience that you need to have the discipline to impose strong abstraction barriers between different modules, or your program will become so complex you can no longer manage nor maintain it.

When you learn a new subject, it helps to have mental abstraction barriers, too, so you know when you can replace one approach with a different approach. I want to give you four levels of abstraction that can help you think about machine learning. It’s important to make mental distinctions between these four things, and the code you write should have modules that reflect these distinctions as well.

<table>
<thead>
<tr>
<th>APPLICATION/DATA</th>
<th>[what kinds of hypotheses are permitted?]</th>
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<tbody>
<tr>
<td>data labeled or not?</td>
<td></td>
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<tr>
<td>yes: labels categorical (classification) or quantitative (regression)?</td>
<td></td>
</tr>
<tr>
<td>no: similarity (clustering) or positioning (dimensionality reduction)?</td>
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<tr>
<th>MODEL</th>
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<td>e.g.:</td>
<td>[what kinds of hypotheses are permitted?]</td>
</tr>
<tr>
<td>– decision fns: linear, polynomial, logistic, neural net, . . .</td>
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<tr>
<td>– nearest neighbors, decision trees</td>
<td></td>
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<tr>
<td>– features</td>
<td></td>
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<tr>
<td>– low vs. high capacity (affects overfitting, underfitting, inference)</td>
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<tr>
<th>OPTIMIZATION PROBLEM</th>
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<td>– variables, objective fn, constraints</td>
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<tr>
<td>e.g., unconstrained, convex program, least squares, PCA</td>
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<table>
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<tr>
<th>OPTIMIZATION ALGORITHM</th>
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<tr>
<td>e.g., gradient descent, simplex, SVD</td>
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In this course, we focus primarily on the middle two levels. As a data scientist, you might be given an application, and your challenge is to turn it into an optimization problem that we know how to solve. We will talk about optimization algorithms, but usually data analysts use optimization codes that are faster and more robust than what they would write themselves.

The second level, the model, has a huge effect on the success of your learning algorithm. Sometimes you get a big improvement by tailoring the model or its features to fit the structure of your specific data. The model also has a big effect on whether you overfit or underfit. And if you want a model that you can interpret so you can do inference, the model has to have a simple structure. Lastly, you have to pick a model that leads to an optimization problem that can be solved. Some optimization problems are just too hard.

It’s important to understand that when you change something in one level of this diagram, you probably have to change all the levels underneath it. If you switch your model from a linear classifier to a neural net, your optimization problem changes, and your optimization algorithm changes too.
[Not all machine learning methods fit this four-level decomposition. Nevertheless, for everything you learn in this class, think about where it fits in this hierarchy. If you don’t distinguish which math is part of the model and which math is part of the optimization algorithm, this course will be very confusing for you.]

**OPTIMIZATION PROBLEMS**

[I want to familiarize you with some types of optimization problems that can be solved reliably and efficiently, and the names of some of the optimization algorithms used to solve them. An important skill for you to develop is to be able to go from an application to a well-defined optimization problem. That skill depends on your ability to recognize well-studied types of optimization problems.]

**Unconstrained**

Goal: Find \( w \) that minimizes (or maximizes) a continuous objective function \( f(w) \).

\( f \) is smooth if its gradient is continuous too.

A **global minimum** of \( f \) is a value \( w \) such that \( f(w) \leq f(v) \) for every \( v \).

A **local minimum** for every \( v \) in a tiny ball centered at \( w \).

[In other words, you cannot walk downhill from \( w \).]

---

```plaintext
minima.pdf
```

Usually, finding a local minimum is easy;
finding the global minimum is hard. [or impossible]

Exception: A function is convex if for every \( x, y \in \mathbb{R}^d \),
the line segment connecting \( (x, f(x)) \) to \( (y, f(y)) \) does not go below \( f(\cdot) \).

---

```plaintext
convex.pdf
```

Formally: for every \( x, y \in \mathbb{R}^d \) and \( \beta \in [0, 1] \), \( f(x + \beta(y - x)) \leq f(x) + \beta(f(y) - f(x)) \).

E.g., perceptron risk fn is convex and nonsmooth.
[When you sum together convex functions, you always get a convex function. The perceptron risk function is a sum of convex loss functions, so it is convex.]

A [continuous] convex function [on a closed, convex domain] has either
- no minimum (goes to $-\infty$), or
- just one local minimum, or
- a connected set of local minima that are all global minima with equal $f$.

[The perceptron risk function is in the last category.]
[In the last two cases, if you walk downhill, you eventually reach a global minimum.]

Gradient descent: repeat $w \leftarrow w - \epsilon \nabla f(w)$

---

[Gradient descent with different learning rates $\epsilon$. Top left: painfully small. Top right: reasonable, but still smaller than ideal. Bottom left: reasonable, but larger than ideal. Bottom right: too large; diverges. This is an animated GIF; see https://gjban.com/machine_learning/gradient_descent_learning_rates.html.]

- Fails/diverges if $\epsilon$ too large.
- Slow if $\epsilon$ too small.
- $\epsilon$ often optimized by trial & error [for slow learners like neural networks].

[The best value of $\epsilon$ is hard to guess. One common technique for dealing with divergence is to check whether a step of gradient descent increases the function value rather than decreasing it; if so, reduce the step size.]

[That’s a simple example of what’s called an adaptive learning rate or a learning rate schedule. These adaptations become even more important when you do stochastic gradient descent or when you optimize non-convex, very twisty objective functions. We’ll revisit the idea when we learn neural networks.]
One interesting aspect of gradient descent that these figures illustrate is that it usually never reaches the exact local minimum. Instead, it gets closer and closer forever, but never exactly reaches the true minimum. We call this behavior “convergence.” The last question of Homework 2 will give you some understanding of why convergence happens under the right conditions.

When we have a feature space with more than one dimension, another problem arises, which is that the learning rate that’s good for one direction might be terrible in another direction. Consider the three examples of gradient descent below.

[Left: 20 iterations of gradient descent on a well-conditioned quadratic function, \( f(w) = 2w_1^2 + w_2^2 \), with a modest step size \( \epsilon = 0.105 \). Center: 20 iterations on an ill-conditioned function, \( f(w) = 10w_1^2 + w_2^2 \); the same step size is now too large. Right: after reducing the step size to \( \epsilon = 0.055 \), we have convergence again but we aren’t approaching the minimum nearly as quickly.]

The step size that works for the left example is too large for the center example; it diverges in the \( w_1 \)-direction. At right, we reduce the step size and obtain convergence. But now convergence is slow in the \( w_2 \)-direction.

High ellipticity of the contours, a.k.a. ill-conditioning of the Hessian, means no learning rate is good in all directions.

The Hessian matrix is said to be ill-conditioned if its largest eigenvalue is much larger than its smallest eigenvalue. Ill-conditioning can be a problem even for simple methods like linear regression, making it harder to solve the problem. In response to these observations, there are adaptive learning rate algorithms that explicitly choose different learning rates in different directions. Famous examples are Adam and RMSprop.

There are many applications where you don’t have a convex objective function. Then gradient descent usually can find a local minimum, but not necessarily a global minimum. And often there is no guarantee that the local minimum you find will be nearly as good as the global minimum. Nevertheless, gradient descent is used for a lot of nonconvex machine learning problems too. For example, neural networks try to optimize an objective function that has lots of local minima. But stochastic gradient descent is still the algorithm of choice for training neural nets. We’ll talk more later in the semester about why.]
**Linear Program**

Linear objective fn + linear inequality constraints.

Goal: Find \( w \) that maximizes (or minimizes) \( c \cdot w \) subject to \( Aw \leq b \)

where \( A \) is \( n \times d \) matrix, \( b \in \mathbb{R}^n \), expressing \( n \) linear constraints:

\[
A_i \cdot w \leq b_i, \quad i \in [1, n]
\]

in \( w \)-space:

---

The set of points \( w \) that satisfy all constraints is a convex polytope called the feasible region \( F \) [shaded].

The optimum is the point in \( F \) that is furthest in the direction \( c \). [What does convex mean?]

A point set \( P \) is convex if for every \( p, q \in P \), the line segment with endpoints \( p, q \) lies entirely in \( P \).

[What is a polytope? Just a polyhedron, generalized to higher dimensions.]

The optimum achieves equality for some constraints (but not most), called the active constraints of the optimum. [In the figure above, there are two active constraints. In an SVM, active constraints correspond to the training points that touch or violate the slab, and these points are also known as support vectors.]

[Sometimes, there is more than one optimal point. For example, in the figure above, if \( c \) pointed straight up, every point on the top horizontal edge would be optimal. The set of optimal points is always convex.]

Example: EVERY feasible point \((w, \alpha)\) gives a linear classifier:

\[
\text{Find } w, \alpha \text{ that satisfies } y_i(w \cdot X_i + \alpha) \geq 1 \text{ for all } i \in [1, n]
\]

[This is the problem of finding a feasible point. This problem can be cast as a slightly different linear program that uses an objective function to make all the inequalities be satisfied strictly if that’s possible.]

IMPORTANT: The data are linearly separable iff the feasible region is not the empty set.

→ Also true for maximum margin classifier (quadratic program)

[The most famous algorithm for linear programming is the simplex algorithm, invented by George Dantzig in 1947. The simplex algorithm is indisputably one of the most important and useful algorithms of the 20th century. It walks along edges of the feasible region, traveling from vertex to vertex until it finds an optimum.]

[Linear programming is very different from unconstrained optimization; it has a much more combinatorial flavor. If you knew which constraints would be the active constraints once you found the solution, it would be easy; the hard part is figuring out which constraints should be the active ones. There are exponentially many possibilities, so you can’t afford to try them all. So linear programming algorithms tend to have a very discrete, computer science feeling to them, like graph algorithms, whereas unconstrained optimization algorithms tend to have a continuous, numerical mathematics feeling.]
Linear programs crop up *everywhere* in engineering and science, but they’re usually in disguise. An extremely useful talent you should develop is to recognize when a problem is a linear program.

A linear program solver can find a linear classifier, but it can’t find the maximum margin classifier. We need something more powerful.

**Quadratic Program**

Quadratic, convex objective fn + linear inequality constraints.

Goal: Find \( w \) that minimizes \[ f(w) = w^T Q w + c^T w \]
subject to \( A w \leq b \)

where \( Q \) is a symmetric, positive semidefinite matrix.

[A matrix is **positive semidefinite** if \( w^T Q w > 0 \) for all \( w \).]

If \( Q \) is positive definite, only one local minimum! [Which is therefore the global minimum.]

[What if \( Q \) is not positive definite? If \( Q \) is indefinite, then \( f \) is not convex, the minimum is not always unique, and quadratic programming is NP-hard. If \( Q \) is positive semidefinite, meaning \( w^T Q w \geq 0 \) for all \( w \), then \( f \) is convex and quadratic programming is tractable, but there may be infinitely many solutions.]

Example: Find maximum margin classifier.

![quadratic.pdf](quadratic.pdf)

Left: A hard-margin SVM minimizes the objective function \( w_1^2 + w_2^2 \). Right: A reminder that there is also an \( \alpha \)-axis, so the isosurfaces of the objective function are really cylinders. On the left isocontours, draw two polygons—one with one active constraint, and one with two—and show the constrained minimum for each polygon. “In a hard-margin SVM, we are looking for the point in this polygon that’s closest to the \( \alpha \)-axis.”

Algs for quadratic programming:

- Simplex-like [commonly used for general-purpose quadratic programs, but not as good for SVMs as the following two algorithms that specifically exploit properties of SVMs]
- Sequential minimal optimization (SMO, used in LIBSVM, “SVC” in scikit)
- Coordinate descent (used in LIBLINEAR, “LinearSVC” in scikit)

Numerical optimization @ Berkeley: EECS 127/227AT/227BT/227C.
6 Decision Theory; Generative and Discriminative Models

DECISION THEORY aka Risk Minimization

[Today I’m going to talk about a style of classifier very different from SVMs. The classifiers we’ll cover in the next few weeks are based on probability.]

[One aspect of probabilistic data is that sometimes a point in feature space doesn’t have just one class. Suppose your data is adult men and women with just one feature: their height. You want to train a classifier that takes in an adult’s height and returns a classification, man or woman. Suppose you are asked to predict the sex of a 5’5” adult. Well, your training set includes some 5’5” women and some 5’5” men. What should you do?]

[In your feature space, you have two training points at the same location with different classes. More generally, the height distributions of men and women overlap. Obviously, in that case, you can’t draw a decision boundary that classifies all points with 100% accuracy.]

Multiple sample points with different classes could lie at same point: we want a probabilistic classifier.

Suppose 10% of population has cancer, 90% doesn’t.

Probability distributions for occupation conditioned on cancer, \(P(X|Y)\):

<table>
<thead>
<tr>
<th>job</th>
<th>(X)</th>
<th>miner</th>
<th>farmer</th>
<th>other</th>
</tr>
</thead>
<tbody>
<tr>
<td>cancer</td>
<td>(Y = 1)</td>
<td>20%</td>
<td>50%</td>
<td>30%</td>
</tr>
<tr>
<td>no cancer</td>
<td>(Y = -1)</td>
<td>1%</td>
<td>10%</td>
<td>89%</td>
</tr>
</tbody>
</table>

[caps here mean random variables, not matrices.]

[I made these numbers up. Please don’t take them as medical advice.]

Recall: \(P(X) = P(X|Y = 1)P(Y = 1) + P(X|Y = -1)P(Y = -1)\)

\[
P(X = \text{farmer}) = 0.5 \times 0.1 + 0.1 \times 0.9 = 0.14 \quad \text{[so 14% of random people are farmers]}
\]

You meet a farmer. Guess whether he has cancer?

[If you’re in a hurry, you might see that 50% of people with cancer are farmers, but only 10% of people with no cancer are farmers, and conclude that a typical farmer probably has cancer. But that would be wrong, because that reasoning fails to take the prior probabilities into account.]

Bayes’ Theorem:

\[
P(Y = 1|X) = \frac{P(X|Y = 1)P(Y = 1)}{P(X)} = \frac{0.05}{0.14}\]

\[
P(Y = -1|X) = \frac{P(X|Y = -1)P(Y = -1)}{P(X)} = \frac{0.09}{0.14} \quad \text{[These two probs always sum to 1.]}\]

\[
P(\text{cancer} | \text{farmer}) = \frac{5}{14} \approx 36%.
\]

[So we probably shouldn’t diagnose cancer.]

[BUT . . . we’re assuming that we want to maximize the chance of a correct prediction. But that’s not always the right assumption. If you’re developing a cheap screening test for cancer, you’d rather have more false positives and fewer false negatives. A false negative might mean somebody misses an early diagnosis and dies of a cancer that could have been treated if caught early. A false positive just means that you spend more money on more accurate tests. When there’s an asymmetry between the awfulness of false positives and false negatives, we can quantify that with a loss function.]
A loss function $L(z, y)$ specifies badness if classifier predicts $z$, true class is $y$.

E.g., $L(z, y) = \begin{cases} 1 & \text{if } z = 1, y = -1, \text{ false positive is bad} \\ 5 & \text{if } z = -1, y = 1, \text{ false negative is BAAAAAD} \\ 0 & \text{if } z = y. \text{ [loss should always be zero for a perfectly correct prediction!]} \end{cases}$

A 36% probability of loss 5 is worse than a 64% prob. of loss 1, so we recommend further cancer screening.

The loss fn above is asymmetrical. [A symmetrical loss is the same for false positives and false negatives. For example . . .]

The 0-1 loss function is $L(z, y) = \begin{cases} 1 & \text{if } z \neq y, \text{ [always 1 for a wrong prediction]} \\ 0 & \text{if } z = y. \text{ [always 0 for a correct prediction]} \end{cases}$

[Another application where you want a very asymmetrical loss function, besides medical diagnosis, is spam detection. Putting a good email in the spam folder is much worse than putting spam in your inbox.]

Let $r : \mathbb{R}^d \to \pm 1$ be a decision rule, aka classifier: a fn that maps a feature vector $x$ to 1 (“in class”) or $-1$ (“not in class”).

**The risk for $r$ is the expected loss over all values of $x, y$:** [Memorize this definition!]

$$R(r) = E[L(r(X), Y)]$$

$$= \sum_x \left( L(r(x), 1) P(Y = 1|X = x) + L(r(x), -1) P(Y = -1|X = x) \right) P(X = x)$$

$$= P(Y = 1) \sum_x L(r(x), 1) P(X = x|Y = 1) + P(Y = -1) \sum_x L(r(x), -1) P(X = x|Y = -1)$$

The Bayes decision rule aka Bayes classifier is the fn $r^*$ that minimizes functional $R(r)$.

Assuming $L(1, 1) = L(-1, -1) = 0$, we have:

$$r^*(x) = \begin{cases} 1 & \text{if } L(-1, 1) P(Y = 1|X = x) > L(1, -1) P(Y = -1|X = x), \\ -1 & \text{otherwise} \end{cases}$$

When $L$ is symmetrical, [the big, key principle you should memorize is]

**pick the class with the biggest posterior probability.**

[But if the loss function is asymmetrical, then you must weight the posteriors with the losses.]

In cancer example, $r^*(miner) = 1$, $r^*(farmer) = 1$, and $r^*(other) = -1$.

The Bayes risk, aka optimal risk, is the risk of the Bayes classifier.

[In our cancer example, the last expression for risk $R$ gives:]

$$R(r^*) = 0.1(5 \times 0.3) + 0.9(1 \times 0.01 + 1 \times 0.1) = 0.249 \quad \text{No decision rule gives a lower risk.}$$

[It is interesting that, if we really know all these probabilities, we really can construct an ideal probabilistic classifier. But in real applications, we rarely know these probabilities; the best we can do is use statistical methods to estimate them.]

Deriving/using $r^*$ is called risk minimization.

[Did you memorize the two boldfaced lines above yet?]
Continuous Distributions

Suppose $X$ has a continuous probability density function (PDF).

Review: [Go back to your CS 70 or stats notes if you don’t remember this.]

\[
\text{prob. that random variable } X \in [x_1, x_2] = \int_{x_1}^{x_2} f(x) \, dx \quad \text{[shaded area]}
\]

area under whole curve $= 1 = \int_{-\infty}^{\infty} f(x) \, dx$

expected value of $g(X) : E[g(X)] = \int_{-\infty}^{\infty} g(x) f(x) \, dx$

mean $\mu = E[X] = \int_{-\infty}^{\infty} x f(x) \, dx$

variance $\sigma^2 = E[(X - \mu)^2] = E[X^2] - \mu^2$

[Perhaps our cancer statistics look like this.]

[Draw this figure by hand (cancerconditional.png)]

[The area under each curve is 1.]

[Let’s use the 0-1 loss function. In other words, suppose you want a classifier that maximizes the chance of a correct prediction. The wrong answer would be to look where these two curves cross and make that be the decision boundary. As before, it’s wrong because it doesn’t take into account the prior probabilities.]

Suppose $P(Y = 1) = 1/3$, $P(Y = -1) = 2/3$, 0-1 loss.

[To maximize the chance you’ll predict correctly whether somebody has cancer, the Bayes decision rule looks up $x$ on this chart and picks the curve with the highest probability. In this example, that means you pick cancer when $x$ is left of the optimal decision boundary, and no cancer when $x$ is to the right.]
Define risk as before, replacing summations with integrals.

\[
R(r) = \mathbb{E}[L(r(X), Y)] = P(Y = 1) \int L(r(x), 1) f_{X|Y=1}(x) \, dx + P(Y = -1) \int L(r(x), -1) f_{X|Y=-1}(x) \, dx.
\]

For Bayes decision rule, Bayes risk is the area under minimum of functions above. [Shade it.]
Assuming \( L(1, 1) = L(-1, -1) = 0 \),

\[
R(r^*) = \int \min_{y=\pm 1} L(-y, y) f_{X|Y=y}(x) \, P(Y = y) \, dx.
\]

[If you want to use an asymmetrical loss function, just scale the curves vertically in the figure above.]

If \( L \) is 0-1 loss, \( R(r) = P(r(x) \text{ is wrong}) \)
and the Bayes optimal decision boundary is \( \{ x : P(Y = 1|X = x) = 0.5 \} \)

\[
\text{Bayes optimal decision boundary}
\]

[Notice that the accuracy of the probabilities is most important near the decision boundary. Far away from the decision boundary, a bit of error in the probabilities probably wouldn’t change the classification.]

[You can also have multi-class classifiers, choosing among three or more classes. The Bayesian approach is a particularly convenient way to generate multi-class classifiers, because you can simply choose whichever class has the greatest posterior probability. Then the decision boundary lies wherever two or more classes are tied for the highest probability.]
3 WAYS TO BUILD CLASSIFIERS

(1) Generative models (e.g., LDA)  [We’ll learn about LDA next lecture.]
   – Assume sample points come from probability distributions, different for each class.
   – Guess form of distributions
   – For each class C, fit distribution parameters to class C points, giving $f_{X|Y=C}(x)$
   – For each C, estimate $P(Y = C)$
   – Bayes’ Theorem gives $P(Y|X)$
     - If 0-1 loss, pick class C that maximizes $P(Y = C|X = x)$ [posterior probability]
       equivalently, maximizes $f_{X|Y=C}(x) P(Y = C)$

(2) Discriminative models (e.g., logistic regression)
   [We’ll learn about logistic regression in a few weeks.]
   – Model $P(Y|X)$ directly

(3) Find decision boundary (e.g., SVM)
   – Model $r(x)$ directly (no posterior)

Advantage of (1 & 2): $P(Y|X)$ tells you probability your guess is wrong
  [This is something SVMs don’t do.]
Advantage of (1): you can diagnose outliers: $f(x)$ is very small
Disadvantages of (1): often hard to estimate distributions accurately;
  real distributions rarely match standard ones.

[What I’ve written here doesn’t actually define the phrases “generative model” or “discriminative model.”
The proper definitions accord with the way statisticians think about models. A generative model is a full
probabilistic model of all variables, whereas a discriminative model provides a model only for the target
variables that we want to predict.]

[It’s important to remember that we rarely know precisely the value of any of these probabilities. There is
usually error in all of these probabilities. In practice, generative models are most popular when you have
phenomena that are well approximated by the normal distribution or another “nice” distribution. Generative
methods also tend to be more stable than other methods when the number of training points is small or when
there are a lot of outliers.]
7 Gaussian Discriminant Analysis, including QDA and LDA

**GAUSSIAN DISCRIMINANT ANALYSIS**

Fundamental assumption: each class has a normal distribution [a Gaussian].

\[ X \sim N(\mu, \sigma^2) : f(x) = \frac{1}{(\sqrt{2\pi}\sigma)^d} \exp\left(-\frac{||x-\mu||^2}{2\sigma^2}\right). \quad [\mu \text{ & } x = \text{vectors}; \sigma = \text{scalar}; d = \text{dimension}]\]

For each class C, suppose we know mean \( \mu_C \) and variance \( \sigma_C^2 \), yielding PDF \( f_{X|Y=C}(x) \), and prior \( \pi_C = P(Y=C) \).

Given \( x \), Bayes decision rule \( r^*(x) \) predicts class C that maximizes \( Q_C(x) = \ln f_{X|Y=C}(x) \pi_C - \ln f_{X|Y=D}(x) \pi_D - d \ln \sigma_C + \ln \pi_C \). [\( Q_C \) is quadratic in \( x \)]

**Quadratic Discriminant Analysis (QDA)**

Suppose only 2 classes C, D. Then the Bayes classifier is

\[
 r^*(x) = \begin{cases} 
 C & \text{if } Q_C(x) - Q_D(x) > 0, \\
 D & \text{otherwise.} 
\end{cases} 
\]

[Picks the class with the biggest posterior probability]
Decision fn is $Q_C(x) - Q_D(x)$ (quadratic); Bayes decision boundary is $Q_C(x) - Q_D(x) = 0$.
- In 1D, B.d.b. may have 1 or 2 points. [Solutions to a quadratic equation]
- In $d$-D, B.d.b. is a quadric. [In 2D, that’s a conic section; see figure above]

[You might not be satisfied with just knowing how each point is classified. One of the great things about QDA is that you can also estimate the probability that your prediction is correct. Let’s work that out.]

To recover posterior probabilities in 2-class case, use Bayes.

$$P(Y = C|X) = \frac{f_{X|Y=C} \pi_C}{f_{X|Y=C} \pi_C + f_{X|Y=D} \pi_D}$$

recall $e^{Q_C(x)} = (\sqrt{2\pi})^d f_{X|Y=C}(x) \pi_C$ [by definition of $Q_C$]

$$P(Y = C|X = x) = \frac{e^{Q_C(x)}}{e^{Q_C(x)} + e^{Q_D(x)}} = \frac{1}{1 + e^{Q_D(x) - Q_C(x)}}$$

$s(\gamma) = \frac{1}{1 + e^{-\gamma}}$ $\Leftarrow$ logistic fn aka sigmoid fn [recall $Q_C - Q_D$ is the decision fn]

[The logistic function. Write beside it:] $s(0) = \frac{1}{2}$, $s(\infty) \to 1$, $s(-\infty) \to 0$, monotonically increasing.

[We interpret $s(0) = \frac{1}{2}$ as saying that on the decision boundary, there’s a 50% chance of class C and a 50% chance of class D.]

Multi-class QDA: [QDA works very naturally with more than 2 classes.]

[Multi-class QDA partitions the feature space into regions. In two or more dimensions, you typically wind up with multiple decision boundaries that adjoin each other at joints. It looks like a sort of Voronoi diagram. In fact, it’s a special kind of Voronoi diagram called a multiplicatively, additively weighted Voronoi diagram.]
Linear Discriminant Analysis (LDA)

[LDA is a variant of QDA with linear decision boundaries. It’s less likely to overfit than QDA.]

Fundamental assumption: all the Gaussians have same variance $\sigma^2$.

[The equations simplify nicely in this case.]

$$Q_C(x) - Q_D(x) = \left(\frac{\mu_C - \mu_D}{\sigma^2}\cdot x\right) - \frac{||\mu_C||^2 - ||\mu_D||^2}{2\sigma^2} + \ln \pi_C - \ln \pi_D.$$ 

[The quadratic terms in $Q_C$ and $Q_D$ canceled each other out!]

Now it’s a linear classifier!

- decision boundary is $w \cdot x + \alpha = 0$
- posterior is $P(Y = C|X = x) = s(w \cdot x + \alpha)$

[The effect of “$s(w \cdot x + \alpha)$” is to scale and translate the logistic fn in $x$-space.]

Special case: if $\pi_C = \pi_D = \frac{1}{2} \Rightarrow (\mu_C - \mu_D) \cdot x - (\mu_C - \mu_D) \cdot \left(\frac{\mu_C + \mu_D}{2}\right) = 0$.

This is the centroid method!

Multi-class LDA: choose $C$ that maximizes linear discriminant fn $\frac{\mu_C \cdot x}{\sigma^2} - \frac{||\mu_C||^2}{2\sigma^2} + \ln \pi_C$.

[Two Gaussians (red) and the logistic function (black). The logistic function is the right Gaussian divided by the sum of the Gaussians. Observe that even when the Gaussians are 2D, the logistic function still looks 1D.]

[Voronoi pdf] [When you have many classes, their LDA decision boundaries form a classical Voronoi diagram if the priors $\pi_C$ are equal. All the Gaussians have the same width.]
MAXIMUM LIKELIHOOD ESTIMATION OF PARAMETERS (Ronald Fisher, circa 1912)

[To use Gaussian discriminant analysis, we must first fit Gaussians to the sample points and estimate the class prior probabilities. We’ll do priors first—they’re easier, because they involve a discrete distribution. Then we’ll fit the Gaussians—they’re less intuitive, because they’re continuous distributions.]

Let’s flip biased coins! Heads with probability $p$; tails w/prob. $1 - p$. [But we don’t know $p$.]

10 flips, 8 heads, 2 tails. [Let me ask you a weird question.] What is the most likely value of $p$?

# of heads is $X \sim \text{B}(n, p)$, binomial distribution:

$P[X = x] = \binom{n}{x} p^x (1 - p)^{n-x}$

[this is the probability of getting exactly $x$ heads in $n$ coin flips]

Prob. of $x = 8$ heads in $n = 10$ flips is

$P[X = 8] = 45p^8 (1 - p)^2 \overset{\text{def}}{=} \mathcal{L}(p)$

Written as a fn of distribution parameter $p$, this prob. is the likelihood fn $\mathcal{L}(p)$.

Maximum likelihood estimation (MLE): A method of estimating the parameters of a statistical model by picking the params that maximize [the likelihood function] $\mathcal{L}$.

. . . is one method of density estimation: estimating a PDF [probability density function] from data.

[Let’s phrase it as an optimization problem.]

Find $p$ that maximizes $\mathcal{L}(p)$.

Solve by finding critical point of $\mathcal{L}$:

$$\frac{d\mathcal{L}}{dp} = 360p^7 (1 - p)^2 - 90p^8 (1 - p) = 0$$

$$\Rightarrow 4(1 - p) - p = 0 \quad \Rightarrow \quad p = 0.8$$

[It shouldn’t seem surprising that a coin that is biased so it comes up heads 80% of the time is the coin most likely to produce 8 heads in 10 flips.]

[Note: $\frac{d^2\mathcal{L}}{dp^2} \overset{\text{def}}{=} -18.9 < 0$ at $p = 0.8$, confirming it’s a maximum.]

[Here’s how this applies to prior probabilities.]

Suppose our training set is $n$ points, with $x$ in class C. Then our estimated prior for class C is $\hat{\pi}_C = x/n$. 
Likelihood of a Gaussian

Given sample points $X_1, X_2, \ldots, X_n$, find best-fit Gaussian.

[Now we want to fit a normal distribution to data, instead of a binomial distribution. If you draw a random point from a normal distribution, what is the probability that it will be exactly at $X_1$?]

[Zero. So it might seem like we have a problem here. With a continuous distribution, the probability of generating any particular point is zero. But we’re just going to ignore that and do “likelihood” anyway.]

Likelihood of drawing these points [in the specified order] is

$$L(\mu, \sigma; X_1, \ldots, X_n) = f(X_1) f(X_2) \cdots f(X_n).$$

[How do we maximize this?]

The log likelihood $\ell(\cdot)$ is the ln of the likelihood $L(\cdot)$.

Maximizing likelihood $\iff$ maximizing log likelihood.

$$\ell(\mu, \sigma; X_1, ..., X_n) = \ln f(X_1) + \ln f(X_2) + \ldots + \ln f(X_n)$$

$$= \sum_{i=1}^{n} \left( - \frac{||X_i - \mu||^2}{2\sigma^2} - d \ln \sqrt{2\pi} - d \ln \sigma \right)$$

In ln of normal PDF

Set $\nabla_\mu \ell = 0$, $\partial \ell / \partial \sigma = 0$ [Find the critical point of $\ell$]

$$\nabla_\mu \ell = \sum_{i=1}^{n} \frac{X_i - \mu}{\sigma^2} = 0 \quad \Rightarrow \quad \hat{\mu} = \frac{1}{n} \sum_{i=1}^{n} X_i$$

[The hats ^ mean “estimated”]

$$\frac{\partial \ell}{\partial \sigma} = \sum_{i=1}^{n} \frac{||X_i - \mu||^2}{\sigma^3} - d \sigma^2 = 0 \quad \Rightarrow \quad \hat{\sigma}^2 = \frac{1}{dn} \sum_{i=1}^{n} ||X_i - \mu||^2$$

We don’t know $\mu$ exactly, so substitute $\hat{\mu}$ for $\mu$ to compute $\hat{\sigma}$.

Takeaway: use sample mean & variance of pts in class C to estimate mean & variance of Gaussian for class C.

For QDA: estimate conditional mean $\hat{\mu}_C$ & conditional variance $\hat{\sigma}_C^2$ of each class C separately [as above] & estimate the priors:

$$\hat{\pi}_C = \frac{n_C}{\sum_D n_D} \iff \text{total sample points in all classes}$$

[$\hat{\pi}_C$ is the coin flip parameter]

For LDA: same means & priors; one variance for all classes:

$$\hat{\sigma}^2 = \frac{1}{dn} \sum_{C} \sum_{(i,y)=C} ||X_i - \hat{\mu}_C||^2 \iff \text{pooled within-class variance}$$

[Notice that although LDA is computing one variance for all the data, each sample point contributes with respect to its own class’s mean. This gives a very different result than if you simply use the global mean! It’s usually smaller than the global variance. We say “within-class” because we use each point’s distance from its class’s mean, but “pooled” because we then pool all the classes together.]
Eigenvectors and the (Anisotropic) Multivariate Normal Distribution

EIGENVECTORS

[I don’t know if you were properly taught about eigenvectors here at Berkeley, but I sure don’t like the way they’re taught in most linear algebra books. So I’ll start with a review. You all know the definition of an eigenvector:]

Given square matrix $A$, if $Av = \lambda v$ for some vector $v \neq 0$, scalar $\lambda$, then $v$ is an eigenvector of $A$ and $\lambda$ is the eigenvalue of $A$ associated w/v.

[But what does that mean? It means that $v$ is a magical vector that, after being multiplied by $A$, still points in the same direction, or in exactly the opposite direction.]

Eigenvalue 2:

Eigenvalue $-\frac{1}{2}$:

[For most matrices, most vectors don’t have this property. So the ones that do are special, and we call them eigenvectors.]
[Clearly, when you scale an eigenvector, it’s still an eigenvector. Only the direction matters, not the length. Let’s look at a few consequences.]

Theorem: if $v$ is eigenvector of $A$ w/eigenvalue $\lambda$, then $v$ is eigenvector of $A^k$ w/eigenvalue $\lambda^k$ [k is a +ve integer; we will use Theorem later]

Proof: $A^2v = A(\lambda v) = \lambda Av = \lambda^2v$, etc.

Theorem: moreover, if $A$ is invertible, then $v$ is eigenvector of $A^{-1}$ w/eigenvalue $1/\lambda$

Proof: $A^{-1}v = A^{-1}(\lambda Av) = \frac{1}{\lambda}v$ [look at the figures above, but go from right to left.]

[Stated simply: When you invert a matrix, the eigenvectors don’t change, but the eigenvalues get inverted. When you square a matrix, the eigenvectors don’t change, but the eigenvalues get squared.]

[Those theorems are pretty obvious. The next theorem is not obvious at all.]
Spectral Theorem: every real, symmetric \( n \times n \) matrix has real eigenvalues and \( n \) eigenvectors that are mutually orthogonal, i.e., \( v_i^\top v_j = 0 \) for all \( i \neq j \)

[This takes about a page of math to prove. One detail is that a matrix can have more than \( n \) eigenvector directions. If two eigenvectors happen to have the same eigenvalue, then every linear combination of those eigenvectors is also an eigenvector. Then you have infinitely many eigenvector directions, but they all span the same plane. So you just arbitrarily pick two vectors in that plane that are orthogonal to each other. By contrast, the set of eigenvalues is always uniquely determined by a matrix, including the multiplicity of the eigenvalues.]

We can use them as a basis for \( \mathbb{R}^n \).

**Building a Matrix with Specified Eigenvectors**

[There are a lot of applications where you’re given a matrix, and you want to extract the eigenvectors and eigenvalues. But when you’re learning the math, I think it’s more intuitive to go in the opposite direction. Suppose you know what eigenvectors and eigenvalues you want, and you want to create the matrix that has those eigenvectors and eigenvalues.]

Choose \( n \) mutually orthogonal unit \( n \)-vectors \( v_1, \ldots, v_n \) [so they specify an orthonormal coordinate system]

Let \( V = [v_1 \ v_2 \ldots \ v_n] \leftarrow n \times n \) matrix

Observe: \( V^\top V = I \) [off-diagonal 0’s because the vectors are orthogonal]

\[
\Rightarrow V^\top = V^{-1} \Rightarrow VV^\top = I
\]

\( V \) is orthonormal matrix: acts like rotation (or reflection)

Choose some eigenvalues \( \lambda_i \):

Let \( \Lambda = \begin{bmatrix} \lambda_1 & 0 & \ldots & 0 \\ 0 & \lambda_2 & 0 \\ \vdots & \ddots & \vdots \\ 0 & 0 & \ldots & \lambda_n \end{bmatrix} \) [diagonal matrix of eigenvalues]

Defn. of eigenvector: \( AV = VA \)

[This is the same definition of eigenvector I gave you at the start of the lecture—\( Av = \lambda v \)—but this version covers all \( n \) eigenvectors in one statement. How do we find the \( A \) that satisfies this equation?]

\[
\Rightarrow AVV^\top = VAV^\top \quad [\text{which proves . . .}]
\]

Theorem: \( A = V\Lambda V^\top = \sum_{i=1}^n \lambda_i \langle v_i | v_i \rangle \) has chosen eigenvectors/values

outer product: \( n \times n \) matrix, rank 1

This is a matrix factorization called the eigendecomposition. [every real, symmetric matrix has one]

Example: [Using the eigenvectors and eigenvalues from the start of the lecture]

\[
A = \begin{bmatrix} 1/\sqrt{2} & 1/\sqrt{2} \\ 1/\sqrt{2} & -1/\sqrt{2} \end{bmatrix} \begin{bmatrix} 2 & 0 \\ 0 & -1/2 \end{bmatrix} \begin{bmatrix} 1/\sqrt{2} & 1/\sqrt{2} \\ 0 & -1/2 \end{bmatrix} = \begin{bmatrix} 3/4 & 5/4 \\ 5/4 & 3/4 \end{bmatrix}.
\]

[This completes our task of finding a symmetric matrix with specified orthonormal eigenvectors and eigenvalues. Again, it is more common in practice that you are given a symmetric matrix, such as a sample covariance matrix, and you need to compute its eigenvectors and eigenvalues. That’s harder. But I think that going from eigenvectors to the matrix helps to build intuition.]
Observe: \( A^2 = V \Lambda V^\top V \Lambda V^\top = V \Lambda^2 V^\top \)
\( A^{-2} = V \Lambda^{-2} V^\top \)
[This is another way to see that squaring a matrix squares its eigenvalues without changing its eigenvectors. It also suggests a way to define a matrix square root.]
Given a symmetric PSD matrix \( \Sigma \), we can find a symmetric square root \( A = \Sigma^{1/2} \):
- compute eigenvectors/values of \( \Sigma \)
- take square roots of \( \Sigma \)'s eigenvalues
- reassemble matrix \( A \) [with the same eigenvectors as \( \Sigma \) but changed eigenvalues]
[Again, the first step of this algorithm—computing the eigenvectors and eigenvalues of a matrix—is much harder than the remaining two steps.]

**Visualizing Quadratic Forms**

[My favorite way to visualize a symmetric matrix is to graph something called the quadratic form, which shows how applying the matrix affects the length of a vector.]
The quadratic form of \( M \) is \( x^\top M x \).
Suppose you want a matrix whose quadratic form has the isocontours at right below, which are circles transformed by \( A \). [The same matrix \( A \) I've been using, which stretches along the direction with eigenvalue 2 and shrinks along the direction with eigenvalue \(-1/2\).]

[Both figures at left are plots of \( \|z\|^2 \), and both figures at right are plots of \( x^\top A^{-2} x \).
(Draw the stretch direction \((1, 1)\) with eigenvalue 2 and the shrink direction \((1, -1)\) with eigenvalue \(-\frac{1}{2}\) on the ellipses at right.)]
That is, we want \( q_2(Az) = q_1(z) \).

Answer: set \( x = Az \).

Then \( q_2(x) = q_1(z) = q_1(A^{-1}x) = \|A^{-1}x\|^2 = x^\top A^{-2}x \).

The isocontours of the quadratic form \( x^\top A^{-2}x \) are ellipsoids determined by the eigenvectors/values of \( A \).

\( \{ x : x^\top A^{-2}x = 1 \} \) is an ellipsoid with axes \( v_1, v_2, \ldots, v_n \) and radii \( \lambda_1, \lambda_2, \ldots, \lambda_n \)

because if \( v_i \) has length 1 (\( v_i \) lies on unit circle), \( x = Av_i \) has length \( \lambda_i \) (\( Av_i \) lies on the ellipsoid).

Therefore, contours of \( x^\top Mx \) are ellipsoids determined by eigenvectors/values of \( M^{-1/2} \).

[The eigenvalues of \( M^{-1/2} \) are the inverse square roots of the eigenvalues of \( M \).]

Special case: \( A \) (or \( M \)) is diagonal \( \IFF \) eigenvectors are coordinate axes
\( \IFF \) ellipsoids are axis-aligned

[Draw axis-aligned isocontours for a diagonal metric.]

A symmetric matrix \( M \) is

<p>| | |</p>
<table>
<thead>
<tr>
<th></th>
<th></th>
</tr>
</thead>
<tbody>
<tr>
<td>positive definite</td>
<td>if ( w^\top Mw &gt; 0 ) for all ( w \neq 0 ) ( \IFF ) all eigenvalues positive</td>
</tr>
<tr>
<td>positive semidefinite</td>
<td>if ( w^\top Mw \geq 0 ) for all ( w ) ( \IFF ) all eigenvalues nonnegative</td>
</tr>
<tr>
<td>indefinite</td>
<td>if +ve eigenvalue &amp; −ve eigenvalue</td>
</tr>
<tr>
<td>invertible</td>
<td>if no zero eigenvalue</td>
</tr>
</tbody>
</table>

[Examples of quadratic forms for positive definite, positive semidefinite, and indefinite matrices. Positive eigenvalues correspond to axes where the curvature goes up; negative eigenvalues correspond to axes where the curvature goes down. (Draw the eigenvector directions, and draw the flat trough in the positive semidefinite bowl.)]

Every squared matrix is pos semidef, including \( A^{-2} \). [Eigenvalues of \( A^{-2} \) are squared, cannot be negative.]

If \( A^{-2} \) exists, it is pos def. [An invertible matrix has no zero eigenvalues.]

What about the isosurfaces of \( x^\top Mx \) for a +ve semidef, \( \text{singular } M \)?

[If \( M \) is only positive semidefinite, but not positive definite, the isosurfaces are cylinders instead of ellipsoids. These cylinders have ellipsoidal cross sections spanning the directions with nonzero eigenvalues, but they run in straight lines along the directions with zero eigenvalues.]
ANISOTROPIC GAUSSIANS

[Let’s revisit the multivariate Gaussian distribution, with different variances along different directions.]

\[ X \sim N(\mu, \Sigma) \quad \quad [X \text{ and } \mu \text{ are } d\text{-vectors. } X \text{ is a random variable with mean } \mu.] \]

\[ f(x) = \frac{1}{\sqrt{(2\pi)^d|\Sigma|}} \exp \left( -\frac{1}{2} (x - \mu)^\top \Sigma^{-1} (x - \mu) \right) \]

↑ determinant of \( \Sigma \)

\( \Sigma \) is the \( d \times d \) SPD covariance matrix.

\( \Sigma^{-1} \) is the \( d \times d \) SPD precision matrix.

Write \( f(x) = n(q(x)) \), where \( q(x) = (x - \mu)^\top \Sigma^{-1} (x - \mu) \)

\[ \mathbb{R} \to \mathbb{R}, \text{ exponential} \quad \mathbb{R}^d \to \mathbb{R}, \text{ quadratic} \]

[Now \( q(x) \) is a function we understand—it’s just a quadratic bowl centered at \( \mu \), the quadratic form of the precision matrix \( \Sigma^{-1} \). The other function \( n(\cdot) \) is a simple, monotonic, convex function, an exponential of the negation of half its argument. This mapping \( n(\cdot) \) does not change the isosurfaces.]

Principle: given monotonic \( n : \mathbb{R} \to \mathbb{R} \), isosurfaces of \( n(q(x)) \) are same as \( q(x) \) (different isovalues).

[(Show this figure on a separate “whiteboard” for easy reuse next lecture.) A paraboloid (left) becomes a bivariate Gaussian (right) after you compose it with a suitable scalar function (center).]
[One of the main ideas is that if you understand the isosurfaces of a quadratic function, then you understand the isosurfaces of a Gaussian, because they’re the same. The differences are in the isovales—in particular, the Gaussian achieves its maximum at the mean, and decreases to zero as you move infinitely far away from the mean.]

The isocontours of \((x - \mu)^\top \Sigma^{-1} (x - \mu)\) are determined by eigenvectors/values of \(\Sigma^{1/2}\).

![Isocontours transformed by \(\Sigma^{1/2}\)]

Aside: \(q(x)\) is the squared distance from \(\Sigma^{-1/2} x\) to \(\Sigma^{-1/2} \mu\). Consider the metric
\[
d(x, \mu) = \left|\left| \Sigma^{-1/2} x - \Sigma^{-1/2} \mu \right|\right| = \sqrt{(x - \mu)^\top \Sigma^{-1} (x - \mu)} = \sqrt{q(x)}.
\]
[So we think of the precision matrix as a “metric tensor” which defines a metric, a sort of warped distance from \(x\) to the mean \(\mu\).]

**Covariance**

Let \(R, S\) be random variables—column vectors or scalars

\[
\text{Cov}(R, S) = \mathbb{E}[(R - \mathbb{E}[R])(S - \mathbb{E}[S])^\top] = \mathbb{E}[RS^\top] - \mu_R \mu_S^\top
\]
\[
\text{Var}(R) = \text{Cov}(R, R)
\]

If \(R\) is a vector, covariance matrix for \(R\) is

\[
\text{Var}(R) = \begin{bmatrix}
\text{Var}(R_1) & \text{Cov}(R_1, R_2) & \ldots & \text{Cov}(R_1, R_d) \\
\text{Cov}(R_2, R_1) & \text{Var}(R_2) & \ldots & \text{Cov}(R_2, R_d) \\
\vdots & \vdots & \ddots & \vdots \\
\text{Cov}(R_d, R_1) & \text{Cov}(R_d, R_2) & \ldots & \text{Var}(R_d) \\
\end{bmatrix}
\]

[Symmetric; each \(R_i\) is scalar]

For a Gaussian \(R \sim \mathcal{N}(\mu, \Sigma)\), one can show \(\text{Var}(R) = \Sigma\). [\ldots as you did in Homework 2.]
[An important point is that statisticians didn’t just arbitrarily decide to call \(\Sigma\) a covariance matrix. Rather, statisticians discovered that if you find the covariance of the normal distribution by integration, it turns out that the covariance is \(\Sigma\). This is a happy fact; it’s rather elegant.]

\(R_i, R_j\) independent \(\Rightarrow\) \(\text{Cov}(R_i, R_j) = 0\) \[the reverse implication is not generally true, but \ldots\]
\(\text{Cov}(R_i, R_j) = 0\) AND multivariate normal dist. \(\Rightarrow\) \(R_i, R_j\) independent

all features pairwise independent \(\Rightarrow\) \(\text{Var}(R)\) is diagonal \[the reverse is not generally true, but \ldots\]

\(\text{Var}(R)\) is diagonal AND joint normal \(\Leftrightarrow\)

\[
f(x) = f(x_1) f(x_2) \cdots f(x_d)
\]

multivariate univariate Gaussians

\(\Rightarrow\) ellipsoids are axis-aligned, with squared radii on diagonal of \(\Sigma = \text{Var}(R)\)

[So when the features are independent, you can write the multivariate Gaussian PDF as a product of univariate Gaussian PDFs. When they aren’t, you can do a change of coordinates to the eigenvector coordinate system, and write it as a product of univariate Gaussian PDFs in eigenvector coordinates. You did something very similar in Q6.2 of Homework 2.]
9 Anisotropic Gaussians: Maximum Likelihood Estimation, QDA, and LDA

ANISOTROPIC GAUSSIANS AND GDA

[Recall from our last lecture the probability density function of the multivariate normal distribution in its full generality. \( x \) and \( \mu \) are \( d \)-vectors.]

Normal PDF: 
\[
f(x) = n(q(x)) = \frac{1}{\sqrt{(2\pi)^d |\Sigma|}} e^{-q/2}, \quad q(x) = (x - \mu)^\top \Sigma^{-1} (x - \mu).
\]

\( \mathbb{R} \rightarrow \mathbb{R} \), exponential \quad \uparrow \quad \text{determinant of } \Sigma \quad \mathbb{R}^d \rightarrow \mathbb{R} \), quadratic

[The covariance matrix \( \Sigma \) and its symmetric square root and its inverse all play roles in our intuition about the multivariate normal distribution. Consider their eigendecompositions.]

\( \Sigma = V\Gamma V^\top \) \quad \text{covariance matrix}
\[\uparrow\] \quad \text{eigenvalues of } \Sigma \text{ are variances along the eigenvectors, } \Gamma_{ii} = \sigma_i^2

\( \Sigma^{1/2} = V\Gamma^{1/2} V^\top \) \quad \text{maps spheres to ellipsoids} \quad (\Sigma^{1/2} \text{ was } A \text{ in last lecture})
\[\uparrow\] \quad \text{eigenvalues of } \Sigma^{1/2} \text{ are Gaussian widths / ellipsoid radii / standard deviations, } \sqrt{\Gamma_{ii}} = \sigma_i

\( \Sigma^{-1} = V\Gamma^{-1} V^\top \) \quad \text{precision matrix (metric tensor)} \quad [\uparrow \text{ quadratic form of } \Sigma^{-1} \text{ defines contours}]

[Recall from last lecture that the isocontours of the multivariate normal distribution are the same as the isocontours of the quadratic form of the precision matrix \( \Sigma^{-1} \).]
Maximum Likelihood Estimation for Anisotropic Gaussians

Given training points \( X_1, \ldots, X_n \) and classes \( y_1, \ldots, y_n \), find best-fit Gaussians.

Let \( n_c = \# \) of training pts in class \( C \).

[Once again, we want to fit the Gaussian that maximizes the likelihood of generating the training points in a specified class. This time I won’t derive the maximum-likelihood Gaussian; I’ll just tell you the answer.]

For QDA:

\[
\hat{\Sigma}_C = \frac{1}{n_c} \sum_{i:y_i = C} (X_i - \hat{\mu}_C)(X_i - \hat{\mu}_C)^T \quad \Leftarrow \text{conditional covariance for pts in class } C
\]

Prior \( \hat{\pi}_C \), mean \( \hat{\mu}_C \): same as last week
[\( \hat{\pi}_C \) is number of points in class \( C \) \div total training points; \( \hat{\mu}_C \) is mean of training points in class \( C \).]

\[\text{maxlike.jpg} \quad \text{[Maximum likelihood estimation takes these points and outputs this Gaussian].} \]

\( \hat{\Sigma}_C \) is positive semidefinite, but not always definite!
[If there are some zero eigenvalues, the standard version of QDA just doesn’t work. We can try to fix it by eliminating the zero-variance dimensions (eigenvectors). Homework 3 suggests a way to do that.]

For LDA:

\[
\hat{\Sigma} = \frac{1}{n} \sum_C \sum_{i:y_i = C} (X_i - \hat{\mu}_C)(X_i - \hat{\mu}_C)^T \quad \Leftarrow \text{pooled within-class covariance matrix}
\]
Let’s revisit QDA and LDA and see what has changed now that we use the multivariate normal distribution in its full, anisotropic generality. The short answer is “not much has changed, but the graphs look cooler.” Conflicting notation warning: capital $X$ represents a random variable, but later it will represent a matrix.

**QDA**

Choosing $C$ that maximizes $P(Y = C|X = x)$ is equivalent to maximizing the quadratic discriminant fn

$$Q_C(x) = \ln \left( \left( \sqrt{2\pi} \right)^d f_{X|Y=C}(x) \pi_C \right) = -\frac{1}{2} (x - \mu_C)^T \Sigma_C^{-1} (x - \mu_C) - \frac{1}{2} \ln |\Sigma_C| + \ln \pi_C.$$  

[This works for any number of classes. In a multi-class problem, you just pick the class with the greatest quadratic discriminant for $x$.]

2 classes: Decision fn $Q_C(x) - Q_D(x)$ is quadratic, but may be indefinite.

$\Rightarrow$ Decision boundary is a quadric.

Posterior is $P(Y = C|X = x) = s(Q_C(x) - Q_D(x))$ where $s(\cdot)$ is logistic fn.

(Show this figure on a separate “whiteboard.”) An example where the decision boundary is a hyperbola—which is not possible with isotropic Gaussians. At left, two anisotropic Gaussians. Center left, the difference $Q_C - Q_D$. After applying the logistic function to this difference we obtain the posterior probabilities at right, which tells us the probability that $x$ is in class $C$. Observe that we can see the decision boundary in both contour plots: it is $Q_C - Q_D = 0$ and $s(Q_C - Q_D) = 0.5$. We don’t need to apply the logistic function to find the decision boundary, but we do need to apply it if we want the posterior probabilities.]
[If we actually know the exact, true parameters $\pi_C$, $\mu_C$, and $\Sigma_C$, this procedure gives us the Bayes optimal decision boundary. By contrast, when we compute approximations $\hat{\pi}_C$, $\hat{\mu}_C$, and $\hat{\Sigma}_C$ from data, that is the QDA algorithm. We hope the QDA decision boundary will approximate the Bayes optimal decision boundary. Sometimes in our textbooks, you will see examples where they plot both the Bayes optimal decision boundary and the decision boundary computed by a learning algorithm. (See the figure two pages forward.) When you see that, the authors know the exact, true probability distributions because they have chosen them and written a program that produces synthetic data from those distributions. With data from the real world, you cannot know the Bayes optimal decision boundary.]

Multi-class QDA:

Multi-class LDA: choose class $C$ that maximizes the linear discriminant fn

$$
\mu_C^\top \Sigma^{-1} x - \frac{\mu_C^\top \Sigma^{-1} \mu_C}{2} + \ln \pi_C.
$$

[Note that we use a linear solver to efficiently compute $\mu_C^\top \Sigma^{-1}$ just once, so the classifier can evaluate test points quickly.]
Anisotropic Gaussians: Maximum Likelihood Estimation, QDA, and LDA

An example of LDA with messy data. The real-world distributions almost surely aren’t Gaussians, but LDA still works reasonably well.

[(Show this figure on a separate “whiteboard.”) In LDA, the decision boundary is always a hyperplane. Note that Mathematica messed up the top left plot a bit; there should be no red in the left corner, nor blue in the right corner.]
Notes on QDA/LDA

- For 2 classes,
  - LDA has $d + 1$ parameters ($w, \alpha$);
  - QDA has $\frac{d(d + 3)}{2} + 1$ params;
  - LDA more likely to underfit;
  - QDA more likely to overfit.  [The danger is much bigger when the dimension $d$ is large.]

[ldaqda.pdf (ISL, Figure 4.9)]  [In these examples, the Bayes optimal decision boundary is purple (and dashed), the QDA decision boundary is green, the LDA decision boundary is black (and dotted). When the Bayes optimal boundary is linear, as at left, LDA gives a more stable fit whereas QDA may overfit. When the Bayes optimal boundary is curved, as at right, QDA often gives you a better fit.]

- QDA on data doesn’t find true optimum Bayes classifier.
  - estimate distributions from finite data.
  - real-world data not perfectly Gaussian.
  - Changing priors or loss = adding constants to discriminant fnns.
    [So it’s very easy. In the 2-class case, it’s equivalent to changing the isovalue … ]
  - Posterior gives decision boundaries for 10% probability, 50%, 90%, etc.
    choosing isovalue = probability $p$ is equivalent to
    - choosing asymmetrical loss $p$ for false positive, $1 - p$ for false negative; OR
    - choosing $\pi_C = 1 - p, \pi_D = p$.
    - With added features, LDA can give nonlinear boundaries; QDA nonquadratic.

[LDA & QDA are the best method in practice for many applications. In the STATLOG project, either LDA or QDA were among the top three classifiers for 10 out of 22 datasets. But it’s not because all those datasets are Gaussian. LDA & QDA work well when the data can only support simple decision boundaries such as linear or quadratic, because Gaussian models provide stable estimates. See ESL, Section 4.3.]
Some Terms

Let $X$ be an $n \times d$ design matrix of sample points. Each row $i$ of $X$ is a sample point $X_i^\top$.

[Now I'm using capital $X$ as a matrix instead of a random variable vector. I'm treating $X_i$ as a column vector to match the standard convention for multivariate PDFs like the Gaussian, but $X_i^\top$ is a row of $X$.]

centering $X$: subtracting $\mu^\top$ from each row of $X$. $X \rightarrow \hat{X}$

[$\mu^\top$ is the mean of all the rows of $X$. Now the mean of all the rows of $\hat{X}$ is zero.]

Let $R$ be drawn from a uniform distribution on sample points. Sample covariance matrix is

$$\text{Var}(R) = \frac{1}{n} \hat{X}^\top \hat{X}.$$  

(This is the simplest way to remember how to compute a covariance matrix for QDA. Imagine you have a design matrix $X_C$ that contains only the sample points of class $C$; then you have $\hat{\Sigma}_C = \frac{1}{n_C} X_C^\top X_C$.)

[When we have points from an anisotropic Gaussian distribution, sometimes it's useful to perform a linear transformation that maps them to an axis-aligned distribution, or maybe even to an isotropic distribution.]

decorrelating $\hat{X}$: applying rotation $Z = \hat{X}V$, where $\text{Var}(R) = V \Lambda V^\top$

[rotates the sample points to the eigenvector coordinate system]

Then $\text{Var}(Z) = \Lambda$. [Z has diagonal covariance. If $X_i \sim \mathcal{N}(\mu, \Sigma)$, then approximately, $Z_i \sim \mathcal{N}(0, \Lambda)$.]

[Proof: $\text{Var}(Z) = \frac{1}{n} Z^\top Z = \frac{1}{n} V^\top \hat{X}^\top \hat{X} V = V^\top \text{Var}(R) V = V^\top \Lambda V^\top V = \Lambda$.]

sphering $\hat{X}$: applying transform $W = \hat{X} \text{Var}(R)^{-1/2}$  [Recall that $\Sigma^{-1/2}$ maps ellipsoids to spheres.]

whitening $X$: centering + sphering, $X \rightarrow W$

Then $W$ has covariance matrix $I$. [If $X_i \sim \mathcal{N}(\mu, \Sigma)$, then approximately, $W_i \sim \mathcal{N}(0, I)$.]

[Whitening input data is often used with other machine learning algorithms, like SVMs and neural networks. The idea is that some features may be much bigger than others—for instance, because they’re measured in different units. SVMs penalize violations by large features more heavily than they penalize small features. Whitening the data before you run an SVM puts the features on an equal basis.]

[One nice thing about discriminant analysis is that whitening is built in.]

[Incidentally, what we’ve done here—computing a sample covariance matrix and its eigenvectors/values—is about 75% of an important unsupervised learning method called principal components analysis, or PCA, which we’ll learn later in the semester.]
10 Regression, including Least-Squares Linear and Logistic Regression

Classification: given point \( x \), predict class (often binary)
Regression: given point \( x \), predict a numerical value

(Classification gives a discrete prediction, whereas regression gives us a quantitative prediction, usually on a continuous scale. We’ve already seen an example of regression in Gaussian discriminant analysis. QDA and LDA don’t just give us a classifier; they also give us the probability that a particular class label is correct. So QDA and LDA do regression on probability values.)

- Choose form of regression fn \( h(x; w) \) with parameters \( w \) (\( h = \) hypothesis)
  - like decision fn in classification [e.g., linear, quadratic, logistic in \( x \)]
- Choose a cost fn (objective fn) to optimize
  - usually based on a loss fn; e.g., empirical risk = expected loss on data

Some regression fns:
(1) linear: \( h(x; w, \alpha) = w \cdot x + \alpha \)
(2) polynomial [equivalent to linear regression with added polynomial features]
(3) logistic: \( h(x; w, \alpha) = s(w \cdot x + \alpha) \) recall: logistic fn \( s(\gamma) = \frac{1}{1+e^{-\gamma}} \)

[The last choice is interesting. You’ll recall that LDA produces a posterior probability function with this expression. So the logistic function seems to be a natural form for modeling certain probabilities. If we want to model posterior probabilities, sometimes we use LDA; but alternatively, we could skip fitting Gaussians to points, and instead just try to directly fit a logistic function to a set of probabilities.]

Some loss fns: let \( z \) be prediction \( h(x) \); \( y \) be true label
(A) \( L(z, y) = (z - y)^2 \) squared error
(B) \( L(z, y) = |z - y| \) absolute error
(C) \( L(z, y) = -y \ln z - (1 - y) \ln(1 - z) \) logistic loss, aka cross-entropy: \( y \in [0,1], z \in (0,1) \)

Some cost fns to minimize:
(a) \( J(h) = \frac{1}{n} \sum_{i=1}^{n} L(h(X_i), y_i) \) mean loss [you can leave out the “\( \frac{1}{n} \)”]
(b) \( J(h) = \max_{i=1}^{n} L(h(X_i), y_i) \) maximum loss
(c) \( J(h) = \sum_{i=1}^{n} \omega_i L(h(X_i), y_i) \) weighted sum [some points are more important than others]
(d) \( J(h) = (a), (b), or (c) + \lambda ||w||^2 \) \( \ell_2 \) penalized/regularized
(e) \( J(h) = (a), (b), or (c) + \lambda ||w||_{\ell_1} \) \( \ell_1 \) penalized/regularized

Some famous regression methods:
Least-squares linear regr.: \( (1) + (A) + (a) \) quadratic cost; minimize w/calculus
Weighted least-sq. linear: \( (1) + (A) + (c) \)
Ridge regression: \( (1) + (A) + (a) + (d) \) quadratic program
Lasso: \( (1) + (A) + (a) + (e) \)
Logistic regr.: \( (3) + (C) + (a) \) convex cost; minimize w/gradient descent
Least absolute deviations: \( (1) + (B) + (a) \) linear program
Chebyshev criterion: \( (1) + (B) + (b) \)

[I have given you several choices of regression function, several choices of loss function, and several choices of objective function. You can snap one part out and replace it with a different one. But the optimization algorithm and its speed depend crucially on which parts you pick. Let’s consider some examples.]
LEAST-SQUARES LINEAR REGRESSION (Gauss, 1801)

Linear regression fn (1) + squared loss fn (A) + cost fn (a).

Find $w, \alpha$ that minimizes $\sum_{i=1}^{n} (X_i \cdot w + \alpha - y_i)^2$.

Convention: $X$ is $n \times d$ design matrix of sample pts
$y$ is $n$-vector of scalar labels

$$
\begin{bmatrix}
X_{11} & X_{12} & \ldots & X_{1j} & \ldots & X_{1d} \\
X_{21} & X_{22} & X_{2j} & X_{2d} \\
\vdots & & & & & \vdots \\
X_{i1} & X_{i2} & X_{ij} & X_{id} \\
\vdots & & & & & \vdots \\
X_{n1} & X_{n2} & X_{nj} & X_{nd}
\end{bmatrix}
\begin{bmatrix}
w_1 \\
w_2 \\
\vdots \\
\alpha
\end{bmatrix}
= 
\begin{bmatrix}
y_1 \\
y_2 \\
\vdots \\
y_n
\end{bmatrix}
$$

[An example of linear regression.]

Usually $n > d$. [But not always.]

Recall fictitious dimension trick [from Lecture 3]: rewrite $h(x) = x \cdot w + \alpha$ as

$$
\begin{bmatrix}
x_1 \\
x_2 \\
1
\end{bmatrix}
\begin{bmatrix}
w_1 \\
w_2 \\
\alpha
\end{bmatrix}
$$

Now $X$ is an $n \times (d + 1)$ matrix; $w$ is a $(d + 1)$-vector. [We’ve added a column of all-1’s to the end of $X$.]

[We rewrite the optimization problem above:]

Find $w$ that minimizes $||Xw - y||^2 = \text{RSS}(w)$, for residual sum of squares
Optimize by calculus:

\[
\text{minimize } \text{RSS}(w) = w^T X^T X w - 2 y^T X w + y^T y \\
\n\n\text{\nabla } \text{RSS} = 2X^T X w - 2X^T y = 0 \\
\Rightarrow \begin{pmatrix} X^T X \end{pmatrix} w = X^T y \\
\text{\(\text{\(d+1\)}} \times \text{\(\text{\(d+1\)}}\text{-vectors} \] \leftarrow \text{the normal equations \([w \text{ unknown; \(X \& y \text{ known}])\]

If \(X^T X\) is singular, problem is underconstrained.
[\ldots because the sample points all lie on a common subspace (through the origin).]
[Notice that \(X^T X\) is always positive semidefinite, but not always positive definite.]

We use a linear solver to find \(w = (X^T X)^{-1} X^T y\). [never actually invert the matrix!]

\(X^+, \text{ the pseudoinverse of } X, (d+1) \times n\)

[We never compute \(X^+\) directly, but we are interested in the fact that \(w\) is a linear transformation of \(y\).]

[X is usually not square, so \(X\) can’t have an inverse. However, every \(X\) has a pseudoinverse \(X^+\), and if \(X^T X\) is invertible, then \(X^+\) is a “left inverse.”]

Observe: \(X^+ X = (X^T X)^{-1} X^T X = I \leftarrow (d + 1) \times (d + 1)\) [which explains the name “left inverse”]

Observe: the predicted values of \(y_i\) are \(\hat{y}_i = w \cdot X_i \Rightarrow \hat{y} = Xw = XX^+y = Hy\)

where \(H = XX^+\) is called the hat matrix because it puts the hat on \(y\).

[Ideally, \(H\) would be the identity matrix and we’d have a perfect fit, but if \(n > d + 1\), then \(H\) is singular.]

Advantages:
- Easy to compute; just solve a linear system.
- Unique, stable solution. [\ldots except when the problem is underconstrained.]

Disadvantages:
- Very sensitive to outliers, because errors are squared!
- Fails if \(X^T X\) is singular. [Which means the problem is underconstrained, has multiple solutions.]

[In discussion section 6, we’ll address how to handle the underconstrained case where \(X^T X\) is singular.]

[Apparently, least-squares linear regression was first posed and solved in 1801 by the great mathematician Carl Friedrich Gauss, who used least-squares regression to predict the trajectory of the planetoid Ceres. A paper he wrote on the topic is regarded as the birth of modern linear algebra.]

**LOGISTIC REGRESSION (David Cox, 1958)**

Logistic regression fn (3) + logistic loss fn (C) + cost fn (a).
Fits “probabilities” in range [0, 1].

Usually used for classification. The input \(y_i\)’s can be probabilities, but in most applications they are all 0 or 1.

QDA, LDA: generative models
logistic regression: discriminative model
[We’ve learned from LDA that in classification, the posterior probabilities are often modeled well by a logistic function. So why not just try to fit a logistic function directly to the data, skipping the Gaussians?]
Regression, including Least-Squares Linear and Logistic Regression

With $X$ and $w$ including the fictitious dimension; $\alpha$ is $w$'s last component . . .

Find $w$ that minimizes

$$J = \sum_{i=1}^{n} L(s(x_i \cdot w), y_i) = -\sum_{i=1}^{n} \left( y_i \ln s(x_i \cdot w) + (1 - y_i) \ln (1 - s(x_i \cdot w)) \right).$$

$J(w)$ is convex! Solve by gradient descent.

[To do gradient descent, we’ll need to compute some derivatives.]

$$s'(\gamma) = \frac{d}{d\gamma} \frac{1}{1 + e^{-\gamma}} = \frac{e^{-\gamma}}{(1 + e^{-\gamma})^2}$$

$$= s(\gamma)(1 - s(\gamma))$$

Let $s_i = s(x_i \cdot w)$

$$\nabla_w J = -\sum \left( \frac{y_i}{s_i} \nabla s_i - \frac{1 - y_i}{1 - s_i} \nabla s_i \right)$$

$$= -\sum \left( \frac{y_i}{s_i} \frac{1 - y_i}{1 - s_i} s_i (1 - s_i) X_i \right)$$

$$= -\sum (y_i - s_i) X_i$$

$$= -X^\top (y - s(Xw))$$

where $s(Xw) = \begin{bmatrix} s_1 \\ s_2 \\ \vdots \\ s_n \end{bmatrix}$

[applies $s$ component-wise to $Xw$]
Gradient descent rule: \( w \leftarrow w + \epsilon X^T (y - s(Xw)) \)

Stochastic gradient descent: \( w \leftarrow w + \epsilon (y_i - s(X_i \cdot w)) X_i \)

Works best if we shuffle points in random order, process one by one.
For very large \( n \), sometimes converges before we visit all points!

[This looks a lot like the perceptron learning rule. The only difference is that the “\(-s_i\)” part is new.]

Starting from \( w = 0 \) works well in practice.

---

If sample pts are linearly separable and \( w \cdot x = 0 \) separates them (with decision boundary touching no pt), scaling \( w \) to have infinite length causes \( s(X_i \cdot w) \rightarrow 1 \) for a pt \( i \) in class \( C \), \( s(X_i \cdot w) \rightarrow 0 \) for a pt not in class \( C \), and \( J(w) \rightarrow 0 \) [in the limit as \( \|w\| \rightarrow \infty \)].

[Moreover, making \( w \) grow extremely large is the only way to get the cost function \( J \) to approach zero.]

Therefore, logistic regression always separates linearly separable pts!

[In this case, the cost function \( J(w) \) has no finite local minimum, but gradient descent will “converge” to a solution, in the sense that the cost \( J \) will get arbitrarily close to zero, though of course the weight vector \( w \) will never become infinitely long. Mathematically speaking, \( w \) doesn’t converge at all—it diverges—though \( J(w) \) does converge to zero.]

[A 2018 paper by Soudry, Hoffer, Nacson, Gunasekar, and Srebro shows that gradient descent applied to logistic regression eventually converges to the maximum margin classifier, but the convergence is very, very slow. A practical logistic regression solver should use a different optimization algorithm.]
11 More Regression; Newton’s Method; ROC Curves

LEAST-SQUARES POLYNOMIAL REGRESSION

Replace each $X_i$ with feature vector $\Phi(X_i)$ with all terms of degree $0 \ldots p$
e.g., $\Phi(X_i) = [X_i^2 \ X_i \ X_i^2 \ X_i \ X_i^2 \ X_i \ X_i^2 \ X_i \ X_i \ X_i]^T$

[Notice that we’ve added the fictitious dimension “1” here, so we don’t need to add it again to do linear or logistic regression. This basis covers all polynomials quadratic in $X_{i1}$ and $X_{i2}$.]

Otherwise just like linear or logistic regression.

Log. reg. + quadratic features = same form of posteriors as QDA.

Very easy to overfit!

[Here are some examples of polynomial overfitting, to show the importance of choosing the polynomial degree very carefully. At left, we have sampled points from a degree-3 curve (black) with added noise. We show best-fit polynomials of degrees 2, 4, 6, and 8 found by regression of the black points. The degree-4 curve (green) fits the true curve (black) well, whereas the degree-2 curve (red) underfits and the degree-6 and 8 curves (blue, yellow) overfit the noise and oscillate. The oscillations in the yellow degree-8 curve are a characteristic problem of polynomial interpolation.]

[At upper right, a degree-20 curve shows just how insane high-degree polynomial oscillations can get. It takes a great deal of densely spaced data to tame the oscillations in a high degree curve, and there isn’t nearly enough data here.]

[At lower right, somebody has regressed a degree-4 curve to U.S. census population numbers. The curve doesn’t oscillate, but can you nevertheless see a flaw? This shows the difficulty of extrapolation outside the range of the data. As a general rule, extrapolation is much harder than interpolation. The $k$-nearest neighbor classifier is one of the few that does extrapolation decently without occasionally returning crazy values.]
This example shows that a fitted degree-10 polynomial (green) can be tamed by using a very large amount of training data (left), even if the training data is noisy. The training data was generated from a different degree-10 polynomial, with noise added. On the right, the same curves are plotted, but the blue diamonds are test points, some of which go outside the range of the training data. We see that the degree-10 regression does decent extrapolation for a short distance, albeit only because the original data was also from a degree-10 polynomial.

**WEIGHTED LEAST-SQUARES REGRESSION**

Linear regression fn (1) + squared loss fn (A) + cost fn (c).

[The idea of weighted least-squares is that some sample points might be more trusted than others, or there might be certain points you want to fit particularly well. So you assign those more trusted points a higher weight. If you suspect some points of being outliers, you can assign them a lower weight.]

Assign each sample pt a weight \( \omega_i \); collect \( \omega_i \)'s in \( n \times n \) diagonal matrix \( \Omega \).

Greater \( \omega_i \) → work harder to minimize \( (\hat{y}_i - y_i)^2 \).[Recall: \( \hat{y}_i \) is predicted label for \( X_i \)]

Find \( w \) that minimizes \( (Xw - y)^\top \Omega (Xw - y) \) = \( \sum_{i=1}^{n} \omega_i (X_i \cdot w - y_i)^2 \)

[As with ordinary least-squares regression, we find the minimum by setting the gradient to zero, which leads us to the normal equations.]

Solve for \( w \) in normal equations: \( X^\top \Omega Xw = X^\top \Omega y \)

**NEWTON'S METHOD**

Iterative optimization method for smooth fn \( J(w) \).

Often much faster than gradient descent. [We’ll use Newton’s method for logistic regression.]

Idea: You’re at point \( v \). Approximate \( J(w) \) near \( v \) by quadratic fn. Jump to its unique critical pt. Repeat until bored.
More Regression; Newton’s Method; ROC Curves

[Three iterations of Newton’s method in one-dimensional space. We seek the minimum of the blue curve, \( J \). Each brown curve is a local quadratic approximation to \( J \). Each iteration, we jump to the bottom of the brown parabola.]

[Steps taken by Newton’s method in two-dimensional space.]

Taylor series about \( v \):

\[
\nabla J(w) = \nabla J(v) + (\nabla^2 J(v)) (w - v) + O(||w - v||^2)
\]

where \( \nabla^2 J(v) \) is the Hessian matrix of \( J \) at \( v \).

Find critical pt \( w \) by setting \( \nabla J(w) = 0 \):

\[
w = v - (\nabla^2 J(v))^{-1} \nabla J(v)
\]

[This is an iterative update rule you can repeat until it converges to a solution. As usual, we probably don’t want to compute a matrix inverse directly. It is faster to solve a linear system of equations, typically by Cholesky factorization or the conjugate gradient method.]

Newton’s method:

1. pick starting point \( w \)
2. repeat until convergence
   - \( e \) ← solution to linear system \((\nabla^2 J(w)) e = -\nabla J(w)\)
   - \( w \) ← \( w + e \)
Warning: Doesn’t know difference between minima, maxima, saddle pts. Starting pt must be “close enough” to desired critical pt.

[If the objective function \( J \) is actually quadratic, Newton’s method needs only one step to find the exact solution. The closer \( J \) is to quadratic, the faster Newton’s method tends to converge.]

[Newton’s method is superior to blind gradient descent for some optimization problems for several reasons. First, it tries to find the right step length to reach the minimum, rather than just walking an arbitrary distance downhill. Second, rather than follow the direction of steepest descent, it tries to choose a better descent direction.]

[Nevertheless, it has some major disadvantages. The biggest one is that computing the Hessian can be quite expensive, and it has to be recomputed every iteration. It can work well for low-dimensional weight spaces, but you would never use it for a neural network, because there are too many weights. Newton’s method also doesn’t work for most nonsmooth functions. It particularly fails for the perceptron risk function, whose Hessian is zero, except where the Hessian is not even defined.]

**LOGISTIC REGRESSION** (continued)

[Let’s use Newton’s method to solve logistic regression faster.]

Recall: \( s'(x) = s(x)(1 - s(x)) \), \( s = [s_1\ s_2\ \ldots\ s_n] \)

\[ \nabla_w J(w) = -\sum_{i=1}^{n} (y_i - s_i) X_i = -X^T(y - s) \]

[Now let’s derive the Hessian too, so we can use Newton’s method.]

\[ \nabla^2_w J(w) = \sum_{i=1}^{n} s_i(1 - s_i) X_i X_i^T = X^T \Omega X \]

where \( \Omega = \begin{bmatrix} s_1(1 - s_1) & 0 & \ldots & 0 \\ 0 & s_2(1 - s_2) & 0 & \ldots \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \ldots & s_n(1 - s_n) \end{bmatrix} \)

\( \Omega \) is +ve definite \( \forall w \Rightarrow X^T \Omega X \) is +ve semidefinite \( \forall w \Rightarrow J \) is convex.

[The logistic regression cost function is convex, so Newton’s method finds a globally optimal point if it converges at all.]

Newton’s method:

\[
\begin{align*}
    w & \leftarrow 0 \\
    & \text{repeat until convergence} \\
    & e \leftarrow \text{solution to normal equations } (X^T \Omega X) e = X^T(y - s) \\
    & \text{Recall: } \Omega, s \text{ are fns of } w \\
    w & \leftarrow w + e
\end{align*}
\]

[Notice that this looks a lot like weighted least squares, but the weight matrix \( \Omega \) and the right-hand-side vector \( y - s \) change every iteration. So we call it . . . ]

An example of iteratively reweighted least squares.
[We need to be very careful with the analogy, though. The weights don’t have the same meaning they had when we learned weighted least-squares regression, because there is no $\Omega$ on the right-hand side of $(X^T \Omega X) e = X^T (y - s)$. Contrary to what you’d expect, a small weight in $\Omega$ causes the Newton iteration to put more emphasis on a point when it computes $e$.]

[Misclassified points far from the decision boundary have the most influence on the step $e$, and correctly classified points far from the decision boundary have the least (because $y_i - s_i$ is small for such a point). Points near the decision boundary have medium influence. But if there are no misclassified points far from the decision boundary, then points near the decision boundary have most of the influence.]

[Here’s one more idea for speeding up logistic regression.] Idea: If $n$ very large, save time by using a random subsample of the pts per iteration. Increase sample size as you go.

[The principle is that the first iteration isn’t going to take you all the way to the optimal point, so why waste time looking at all the sample points? Whereas the last iteration should be the most accurate one.]

**LDA vs. Logistic Regression**

Advantages of LDA:
- For well-separated classes, LDA stable; log. reg. surprisingly unstable
- > 2 classes easy & elegant; log. reg. needs modifying (softmax regression) [see Discussion 6]
- LDA slightly more accurate when classes nearly normal, especially if $n$ is small

Advantages of log. reg.:
- More emphasis on decision boundary; always separates linearly separable pts
  [Correctly classified points far from the decision boundary have a small effect on logistic regression—albeit a bigger effect than they have on SVMs—whereas misclassified points far from the decision boundary have the biggest effect. By contrast, LDA gives all the sample points equal weight when fitting Gaussians to them. Weighting points according to how badly they’re misclassified is good for reducing training error, but it can also be bad if you want stability or insensitivity to bad data.]

- More robust on some non-Gaussian distributions (e.g., dists. w/large skew)
- Naturally fits labels between 0 and 1 [usually probabilities]

[When you use logistic regression with added quadratic features, you get a quadric decision boundary, just as you do with QDA. Based on what I’ve said here, do you think logistic regression with quadratic features gives you exactly the same classifier as QDA?]
**ROC CURVES** (for test sets)

[This is a ROC curve. That stands for receiver operating characteristics, which is an awful name but we’re stuck with it for historical reasons.

A ROC curve is a way to evaluate your classifier after it is trained. It is made by running a classifier on the test set or validation set. It shows the rate of false positives vs. true positives for a range of settings.

We assume there is a knob we can turn to trade off these two types of error. For our purposes, that knob is the posterior probability threshold for Gaussian discriminant analysis or logistic regression. However, neither axis of this plot is that knob.]

*x-axis: “false positive rate = % of −ve classified as +ve”*

*y-axis: “true positive rate = % of +ve classified as +ve aka sensitivity”*

“false negative rate”: vertical distance from curve to top [1 − sensitivity]

“specificity”: horizontal distance from curve to right [1 − false positive rate; “true negative rate”]

[You generate this curve by trying every probability threshold; for each threshold, measure the false positive & true positive rates and plot a point.]

upper right corner: “always classify +ve (Pr ≥ 0)”

lower left corner: “always classify −ve (Pr > 1)”

diagonal: “random classifiers”

[A rough measure of a classifier’s effectiveness is the area under the curve. For a classifier that is always correct, the area under the curve is one. For the random classifier, the area under the curve is 1/2, so you’d better do better than that.]

IMPORTANT: In practice, the trade-off between false negatives and false positives is usually negotiated by choosing a point on this plot, based on real test data, and NOT by taking the choice of threshold that’s best in theory.
12 Statistical Justifications; the Bias-Variance Decomposition

STATISTICAL JUSTIFICATIONS FOR REGRESSION

[So far, I’ve talked about regression as a way to fit curves to points. Recall that early in the semester I divided machine learning into 4 levels: the application, the model, the optimization problem, and the optimization algorithm. My last two lectures about regression were at the bottom two levels: optimization. But why did we pick these cost functions? Today, let’s take a step up to the second level, the model. I will describe some models, how they lead to those optimization problems, and how they contribute to underfitting or overfitting.]

Typical model of reality:
- sample points come from unknown prob. distribution: $X_i \sim D$.
- $y$-values are sum of unknown, non-random fn + random noise: $\forall X_i, \ y_i = g(X_i) + \epsilon_i, \ \epsilon_i \sim D'$, $D'$ has mean zero.

[We are positing that reality is described by a function $g$. We don’t know $g$, but $g$ is not random; it represents a consistent relationship between $X$ and $y$ that we can estimate. We add to $g$ a random variable $\epsilon$, which represents measurement errors and all the other sources of statistical error when we measure real-world phenomena. Notice that the noise is independent of $X$. That’s a pretty questionable assumption, and often it does not apply in practice, but that’s all we’ll have time to deal with this semester. Also notice that this model leaves out systematic errors, like when your measuring device adds one to every measurement, because we usually can’t diagnose systematic errors from data alone.]

Goal of regression: find $h$ that estimates $g$.
Ideal approach: choose $h(x) = E_Y[Y|X=x] = g(x) + E[\epsilon] = g(x)$.

[If this expectation exists at all, it partly justifies our model of reality. We can retroactively define $g$ to be this expectation.]

Least-Squares Cost Function from Maximum Likelihood

Suppose $\epsilon_i \sim N(0, \sigma^2)$; then $y_i \sim N(g(X_i), \sigma^2)$.

Recall that log of normal PDF is
$$\ln f(y_i) = -\frac{(y_i - \mu)^2}{2\sigma^2} - \text{constant} \quad \Leftarrow \mu_i = g(X_i)$$

& log likelihood is
$$\ell(g; X, y) = \ln (f(y_1) f(y_2) \cdots f(y_n)) = \ln f(y_1) + \ldots + \ln f(y_n) = -\frac{1}{2\sigma^2} \sum (y_i - g(X_i))^2 - \text{constant.}$$

Takeaway: Max likelihood on “parameter” $g \Rightarrow$ choose a $g$ that minimizes $\sum(y_i - g(X_i))^2$.

[We treat $g$ as a “distribution parameter.” If the noise is normally distributed, maximum likelihood tells us to estimate $g$ by least-squares regression.]

[However, I’ve told you in previous lectures that least-squares is very sensitive to outliers. If the error is truly normally distributed, that’s not a big deal, especially when you have a lot of sample points. But in the real world, the distribution of outliers often isn’t normal. Outliers might come from wrongly measured measurements, data entry errors, anomalous events, or just not having a normal distribution. When you have a heavy-tailed distribution of noise, for example, least-squares isn’t a good choice.]
Empirical Risk

The risk for hypothesis $h$ is expected loss $R(h) = E[L]$ over all $(X, Y)$ in some joint distribution. Discriminative model: we don’t know $X$’s dist. $D$. How can we minimize risk?

[If we have a generative model, we can estimate the joint probability distribution for $X$ and $Y$ and derive the expected loss. That’s what we did for Gaussian discriminant analysis. But today I’m assuming we don’t have a generative model, so we don’t know those probabilities. Instead, we approximate the distribution in a very crude way: we pretend that the sample points are the distribution.]

Empirical distribution: the discrete uniform distribution over the sample pts

Empirical risk: expected loss under empirical distribution

$$\hat{R}(h) = \frac{1}{n} \sum_{i=1}^{n} L(h(X_i), y_i)$$

[The hat on the $R$ indicates that $\hat{R}$ is only a cheap approximation of the true, unknown statistical risk $R$ we really want to minimize. Often, this is the best we can do. For many but not all distributions, the empirical risk converges to the true risk in the limit as $n \to \infty$. Choosing $h$ that minimizes $\hat{R}$ is called empirical risk minimization.]

Takeaway: this is why we [usually] minimize the average of the loss fns.

Logistic Loss from Maximum Likelihood

What cost fn should we use for probabilities?

Actual probability pt $X_i$ is in class C is $y_i$; predicted prob. is $h(X_i)$.

Imagine $\beta$ duplicate copies of $X_i$, with $y_i\beta$ in class C, and $(1 - y_i)\beta$ not.

[The size of $\beta$ isn’t very important, but imagine that $y_i\beta$ and $(1 - y_i)\beta$ are both integers for all $i$. If $y_i$ is irrational, approximate it with a very close rational number.]

[Let’s use maximum likelihood estimation to choose the hypothesis most likely to generate these labels for these sample points. The following likelihood is the probability of generating these labels in a particular fixed order.]

Likelihood is $L(h; X, y) = \prod_{i=1}^{n} h(X_i)^{y_i\beta}(1 - h(X_i))^{(1-y_i)\beta}$

Log likelihood is $\ell(h) = \ln L(h)$

$$= \beta \sum_i \left( y_i \ln h(X_i) + (1 - y_i) \ln(1 - h(X_i)) \right)$$

$$= -\beta \sum \text{logistic loss fn } L(h(X_i), y_i).$$

Takeaway: Max likelihood $\Rightarrow$ minimize $\sum$ logistic losses.

[So the principle of maximum likelihood explains where the weird logistic loss function comes from.]
THE BIAS-VARIANCE DECOMPOSITION

There are 2 sources of error in a hypothesis $h$:

*bias*: error due to inability of hypothesis $h$ to fit $g$ perfectly
e.g., fitting quadratic $g$ with a linear $h$

*variance*: error due to fitting random noise in data
e.g., we fit linear $g$ with a linear $h$, yet $h \neq g$.

Model: \[X_i \sim D, \epsilon_i \sim D', y_i = g(X_i) + \epsilon_i\] \[\text{[remember that } D'\text{ has mean zero]}\]
fit hypothesis $h$ to $X, y$

Now $h$ is a random variable; i.e., its weights are random

Consider arbitrary pt $z \in \mathbb{R}^d$ (not necessarily a sample pt!) & $\gamma = g(z) + \epsilon$, $\epsilon \sim D'$
[So $z$ is arbitrary, whereas $\gamma$ is random.]

Note: $E[\gamma] = g(z)$; $\text{Var}(\gamma) = \text{Var}(\epsilon)$ \[\text{[the mean comes from } g, \text{ and the variance comes from } \epsilon]\]

Risk fn when loss = squared error:

\[R(h) = E[L(h(z), \gamma)]\]
\[\uparrow \text{take expectation over possible training sets } X, y \& \text{values of } \gamma\]
[Stop and take a close look at this expectation. Remember that the hypothesis $h$ is a random variable. We are taking a mean over the probability distribution of hypotheses. That seems pretty weird if you’ve never seen it before. But remember, the training data $X$ and $y$ come from a joint probability distribution. We use the training data to choose weights, so the weights that define $h$ also come from some probability distribution. It might be hard to work out what that distribution is, but it exists. This “$E[\cdot]$” is integrating the loss over all possible values of the weights.]

\[= E[(h(z) - \gamma)^2]\]
\[= E[h(z)^2] + E[\gamma^2] - 2 \text{E}[\gamma h(z)]\]
\[\text{[Observe that } \gamma \text{ and } h(z) \text{ are independent]}\]
\[= \text{Var}(h(z)) + E[h(z)^2] + \text{Var}(\gamma) + E[\gamma]^2 - 2E[\gamma] E[h(z)]\]
\[= E[h(z)^2] - E[\gamma]^2 + \text{Var}(h(z)) + \text{Var}(\gamma)\]
\[= \frac{(E[h(z)] - g(z))^2}{\text{bias}^2 \text{ of method}} + \frac{\text{Var}(h(z))}{\text{variance of method}} + \frac{\text{Var}(\epsilon)}{\text{irreducible error}}\]

[This is called the *bias-variance decomposition* of the risk function. Let’s look at an intuitive interpretation of these three parts.]
In this example, we’re trying to fit a sine wave with lines, which obviously aren’t going to be accurate. At left, we have generated 50 different hypotheses (lines). Each line was generated from 20 random training points by least-squares linear regression. At upper right, the red line is the expected hypothesis—an average over infinitely many hypotheses. The black curve illustrates test points on the true function \( g \). We see that most test points have a large bias (difference between the black and red curves), because lines don’t fit sine waves well. However, some of the test points happen to have a small bias—where the sine wave crosses the red line. At center right, the variance is the expected squared difference between a random black line and the red line (at a test point \( z \)). At lower right, the irreducible error is the expected squared difference between a random test point and the sine wave.

This is pointwise version [of the bias-variance decomposition.]

Mean version: let \( z \sim D \) be random variable; take mean over \( D \) of \( \text{bias}^2 \), variance.

[So you can decompose one test point’s error into these three numbers, or you can decompose the error of the hypothesis over its entire range into three numbers, which tells you roughly how big they’ll be on a large test set.]
[Now I will write down a list of consequences of what we’ve just learned.]

- Underfitting = too much bias
- Most overfitting caused by too much variance
- Training error reflects bias but not much variance; test error reflects both
  [which is why low training error can fool you when you’ve overfitted]
- For many distributions, variance $\to 0$ as $n \to \infty$
- If $h$ can fit $g$ exactly, for many distributions bias $\to 0$ as $n \to \infty$
- If $h$ cannot fit $g$ well, bias is large at “most” points
- Adding a good feature reduces bias; adding a bad feature rarely increases it
- Adding a feature usually increases variance  [don’t add a feature unless it reduces bias more]
- Can’t reduce irreducible error [hence its name]
- Noise in test set affects only Var($\epsilon$);
  noise in training set affects only bias & Var($h$)
- We can’t precisely measure bias or variance of real-world data
  [because we cannot know $g$ exactly and our noise model might be wrong]
- But we can test learning alg’s by choosing $g$ & making synthetic data

[At left, a data set is fit with splines having various degrees of freedom. The synthetic data is taken from the black curve with added noise. At center, we plot training error (gray) and test error (red) as a function of the number of degrees of freedom. At right, we plot the squared test error as a sum of squared bias (blue) and variance (orange). As the number of degrees of freedom increases, the training and test errors both decrease up to 6 degrees because the bias decreases, but for more degrees of freedom the test error increases because the variance increases.]
Example: Least-Squares Linear Reg.

For simplicity, no fictitious dimension.
[This implies that our linear regression function has to be zero at the origin.]

Model: \( g(z) = v^T z \) (ground truth is linear)
[So we could fit \( g \) perfectly with a linear \( h \) if not for the noise in the training set.]
Let \( e \) be noise \( n \)-vector, \( e_i \sim \mathcal{N}(0, \sigma^2) \)
Training labels: \( y = Xv + e \)
[X & \( y \) are the inputs to linear regression. We don’t know \( v \) or \( e \).]

Lin. reg. computes weights
\[
w = X^+ y = X^+(Xv + e) = v + X^+ e
\]
[We want \( w = v \), but the noise in \( y \) becomes noise in \( w \).]

noise in weights

BIAS is \( \text{E}[h(x)] - g(x) = \text{E}[w^T z] - v^T z = z^T \text{E}[w - v] = z^T \text{E}[Xv] = z^T \text{E}[X^+] \text{E}[e] = 0 \)

Warning: This does not mean \( h(x) - g(x) \) is always 0!
[Those deviations from the mean are captured in the variance.]
When the bias is zero, a perfect fit is possible. But when a perfect fit is possible, not all learning methods
give you a bias of zero; here it’s a benefit of the squared error loss function. With a different noise or a
different loss function, we might have a nonzero bias even fitting a linear \( h \) to a linear \( g \).]

VARIANCE is \( \text{Var}(h(x)) = \text{Var}(w^T z) = \text{Var}(v^T z + (X^+ e)^T z) = \text{Var}(z^T X^+ e) \)
[This is the dot product of a vector \( z^T X^+ \) with an isotropic, normally distributed vector \( e \). The dot product
reduces it to a one-dimensional Gaussian along the direction \( z^T X^+ \), so this variance is just the variance of
the 1D Gaussian times the squared length of the vector \( z^T X^+ \).]
\[
\begin{align*}
\text{Var}(h(x)) &= \sigma^2 \left\| z^T X^+ \right\|^2 = \sigma^2 z^T (X^T X)^{-1} X^T X (X^T X)^{-1} z \\
&= \sigma^2 z^T (X^T X)^{-1} z
\end{align*}
\]
If we choose coordinate system so \( D \) has mean zero, then \( X^T X \to n \text{Var}(D) \) as \( n \to \infty \), so for \( z \sim D \),
\[
\text{Var}(h(x)) \approx \sigma^2 \frac{d}{n}
\]
[where \( d \) is the dimension—the number of features per sample point.]

[Why? With the eigendecomposition \( \text{Var}(D) = \Lambda V^T \), we have \( \text{E}[z^T \text{Var}(D)^{-1} z] = \text{E}[\|\Lambda^{-1/2} V^T z\|^2] = \sum_{i=1}^d \text{E}[(v_i \cdot z)^2]/\Lambda_i. \) But as \( z \sim D, \text{E}[(v_i \cdot z)^2] = \text{Var}[v_i \cdot z] = \Lambda_i, \) so \( \text{E}[z^T \text{Var}(D)^{-1} z] = d. \) Approximating the
covariance \( \text{Var}(D) \) with the sample covariance matrix gives \( \text{E}[z^T (X^T X)^{-1} z] \approx d/n. \)]

Takeaways: Bias can be zero when hypothesis function can fit the real one!
[This is a nice property of the squared error loss function.]
Variance portion of RSS (overfitting) decreases as \( 1/n \) (sample points),
\[
\text{increases as } d \text{ (features)}
or \( O(d^p) \) if you use degree-\( p \) polynomials.
\]
[I’ve used linear regression because it’s a relatively simple example. But the bias-variance trade-off applies
to many learning algorithms, including classification as well as regression. But for most learning algorithms,
the math gets a lot more complicated than this, if you can do it at all. Sometimes there are multiple competing
bias-variance models and no consensus on which is the right one.]
13 Shrinkage: Ridge Regression, Subset Selection, and Lasso

RIDGE REGRESSION aka Tikhonov Regularization

Least-squares linear regression + $\ell_2$ penalized mean loss. (1) + (A) + (a) + (d).

Find $w$ that minimizes $\|Xw - y\|^2 + \lambda \|w'\|^2 = J(w)$
where $w'$ is $w$ with component $\alpha$ replaced by 0.

$X$ has fictitious dimension but we DON’T penalize $\alpha$.

Adds a regularization term, aka a penalty term, for shrinkage: to encourage small $\|w'\|$. Why?

- Guarantees positive definite normal eq’ns; always unique solution.
  [Standard least-squares linear regression yields singular normal equations when the sample points lie on a common hyperplane in feature space—for example, when $d > n$.]

[The cost function $J(w)$ without and with regularization. This plot ignores the dimension of the bias term $\alpha$.]

[At left, we see a cost function for least-squares regression, a positive semidefinite quadratic form. This cost function has many minima, and the regression problem is said to be ill-posed. By adding a small penalty term, we obtain a positive definite quadratic form (right), with one unique minimum. “Regularization” implies that we are turning an ill-posed problem into a well-posed problem.]

[That was the original motivation, but the next has become more important in machine learning . . .]

- Reduces overfitting by reducing variance. Why?

  Example: Input $X_1 = (0, 0)$ with label 0; $X_2 = (1, 1)$ with label 0; $X_3 = (0.51, 0.49)$ with label 1.
  Linear regr. gives $50x_1 - 50x_2$. [This linear function fits all three points exactly.]
  Big weights!

  [Weights this big would be justified if there were big differences between labels, or if there were small distances between points, but neither is true. Large weights imply that tiny changes in $x$ can cause huge changes in $y$. Consider that the labels don’t differ by more than 1 and the points are separated by distances greater than 0.7. So these disproportionately large weights are a sure sign of overfitting.] So we penalize large weights.

  [This use of regularization is closely related to the first one. When you have large variance and a lot of overfitting, it implies that your problem is close to being ill-posed, even though technically it might be well-posed.]
Jonathan Richard Shewchuk

isocontours of \( \|Xw - y\|_2^2 \)

isocontours of \( \|w\|_2^2 \)

least-squares solution

ridge solution for several values of \( \lambda \)

[ridgeterms2.pdf (redrawing of ISL, Figure 6.7)] In this plot of weight space, for simplicity, we’re not using a bias term \( \alpha \) (we set it to zero). \( \hat{w} \) is the least-squares solution. The red ellipses are isocontours of \( \|Xw - y\|_2^2 \). The blue circles are isocontours of \( \|w\|_2^2 \), centered at the origin. The ridge regression solution lies where a red isocontour just touches a blue isocontour tangentially. As \( \lambda \) increases, the solution will occur at a more outer red isocontour and a more inner blue isocontour. This shrinks \( w \) and helps to reduce overfitting.

Setting \( \nabla J = 0 \) gives normal eq’ns

\[
(X^T X + \lambda I') \hat{w} = X^T y
\]

where \( I' \) is identity matrix w/bottom right set to zero. [Don’t penalize the bias term \( \alpha \).]

[Don’t worry; \( X^T X + \lambda I' \) is always positive definite for \( \lambda > 0 \), assuming \( X \) ends with a column of 1’s.]

Algorithm: Solve for \( \hat{w} \). Return \( h(z) = \hat{w}^T z \).

Increasing \( \lambda \) ⇒ more regularization; smaller \( \|w'\| \)

Recall [from the previous lecture] our data model \( y = Xv + e \), where \( e \) is noise.

Variance of ridge regr. at test pt \( z \) is \( \text{Var}(z^T (X^T X + \lambda I')^{-1} X^T e) \).

As \( \lambda \to \infty \), variance \( \to 0 \), but bias increases.

[ridgebiasvar.pdf (ISL, Figure 6.5)] [Plot of bias\(^2 \) & variance as \( \lambda \) increases.]

[So, as usual for the bias-variance trade-off, the test error as a function of \( \lambda \) is a U-shaped curve. We find the bottom by validation.]

\( \lambda \) is a hyperparameter; tune by (cross-)validation.

Ideally, features should be “normalized” to have same variance.

Alternative: use asymmetric penalty by replacing \( I' \) w/other diagonal matrix. [For example, if you use polynomial features, you could use different penalties for monomials of different degrees.]
Bayesian Justification for Ridge Reg.

Assign a prior probability on \( w' \): \( w' \sim N(0, \varsigma^2) \), with PDF \( f(w') \propto e^{-\|w'\|^2/(2\varsigma^2)} \)

[This prior probability says that we think weights close to zero are more likely to be correct.]

Apply MLE to maximize the posterior prob.

Bayes’ Theorem: posterior \( f_{W|X,Y}(w) = \frac{f_{Y|X,W}(y) f(w')}{f_{Y|X}(y)} \)

Maximize log posterior = \( \ln f_{Y|X,W}(y) + \ln f(w') - \text{const} \)

= \( -\text{const} \|Xw - y\|^2 - \text{const} \|w'\|^2 - \text{const} \)

\( \Rightarrow \) Minimize \( \|Xw - y\|^2 + \lambda \|w'\|^2 \)

[We are treating \( w \) and \( y \) as random variables, but \( X \) as a fixed constant—it’s not random.]

This method (using MLE, but maximizing posterior) is called maximum \textit{a posteriori} (MAP).

[A prior probability on the weights is another way to understand regularizing ill-posed problems.]

FEATURE SUBSET SELECTION

[Some of you may have noticed as early as Homework 1 that you can sometimes get better performance on a spam classifier simply by dropping some useless features.]

All features increase variance, but not all features reduce bias.

Idea: Identify poorly predictive features, ignore them (weight zero).

Less overfitting, smaller test errors.

2nd motivation: Inference. Simpler models convey interpretable wisdom.

Useful in all classification & regression methods.

Sometimes it’s hard: Different features can partly encode same information.

Combinatorially hard to choose best feature subset.

Alg: Best subset selection. Try all \( 2^d - 1 \) nonempty subsets of features. [Train one classifier per subset.]

Choose best classifier by (cross-)validation. Slow.

[Obviously, best subset selection isn’t feasible if we have a lot of features. But it gives us an “ideal” algorithm to compare practical algorithms with. If \( d \) is large, there is no algorithm that’s guaranteed to find the best subset and that runs in acceptable time. But heuristics often work well.]

Heuristic 1: Forward stepwise selection.

Start with \textbf{null model} (0 features); repeatedly add best feature until validation errors start increasing (due to overfitting) instead of decreasing. At each outer iteration, inner loop tries every feature & chooses the best by validation. Requires training \( O(d^2) \) models instead of \( O(2^d) \).

Not perfect: e.g., won’t find the best 2-feature model if neither of those features yields the best 1-feature model. [That’s why it’s a heuristic.]

Heuristic 2: Backward stepwise selection.

Start with all \( d \) features; repeatedly remove feature whose removal gives best reduction in validation error. Also trains \( O(d^2) \) models.

[Forward stepwise is a better choice when you suspect only a few features will be good predictors; e.g., spam. Backward stepwise is better when most features are important. If you’re lucky, you’ll stop early.]
**LASSO (Robert Tibshirani, 1996)**

Least-squares linear regression + $\ell_1$ penalized mean loss. (1) + (A) + (a) + (e).

“Least absolute shrinkage and selection operator”

[This is a regularized regression method similar to ridge regression, but it has the advantage that it often naturally sets some of the weights to zero.]

Find $w$ that minimizes $\|Xw - y\|^2 + \lambda \|w\|_1$ where $\|w\|_1 = \sum_{i=1}^{d} |w_i|$ (Don’t penalize $\alpha$.)

Recall ridge regr.: isosurfaces of $\|w\|^2$ are hyperspheres.

The isosurfaces of $\|w\|_1$ are cross-polytopes.

The unit cross-polytope is the convex hull of all the positive & negative unit coordinate vectors.

\[ \|w\|_1 = 1 \]

[Draw this figure by hand](crosspolys.png)

[You get larger and smaller cross-polytope isosurfaces by scaling these.]

\[ \text{Isocontours of the terms of the objective function for the Lasso appear at left. Compare with the ridge regression isocontours at right.} \]

[The red ellipses are the isocontours of $\|Xw - y\|^2$, and the least-squares solution lies at their center. The isocontours of $\|w\|_1$ are diamonds centered at the origin (blue). The solution lies where a red isocontour just touches a blue diamond. What’s interesting here is that in this example, the red isocontour touches just the tip of the diamond. So the weight $w_1$ gets set to zero. That’s what we want to happen to weights that don’t have enough influence. This doesn’t always happen—for instance, the red isosurface could touch a side of the diamond instead of a tip of the diamond.]

[When you go to higher dimensions, you might have several weights set to zero. For example, in 3D, if the red isosurface touches a sharp vertex of the cross-polytope, two of the three weights get set to zero. If it touches a sharp edge of the cross-polytope, one weight gets set to zero. If it touches a flat side of the cross-polytope, no weight is zero.]
[This shows the weights for a typical linear regression problem with about 10 variables. You can see that as lambda increases, more and more of the weights become zero. Only four of the weights are really useful for prediction; they’re in color. Statisticians used to choose $\lambda$ by looking at a chart like this and trying to eyeball a spot where there aren’t too many predictors and the weights aren’t changing too fast. But nowadays they prefer validation.]

Sometimes sets some weights to zero, especially for large $\lambda$.

Algs: subgradient descent, least-angle regression (LARS), forward stagewise

[Lasso can be reformulated as a quadratic program, but it’s a quadratic program with $2^d$ constraints, because a $d$-dimensional cross-polytope has $2^d$ facets. In practice, special-purpose optimization methods have been developed for Lasso. I’m not going to teach you one, but if you need one, look up the last two of these algorithms. LARS is built into the R Programming Language for statistics.]

[As with ridge regression, you should probably normalize the features first before applying Lasso.]
14 Decision Trees

**DECISION TREES**

Nonlinear method for classification and regression.

Uses tree with 2 node types:
- internal nodes test feature values (usually just one) & branch accordingly
- leaf nodes specify class \( h(x) \)

![Decision Tree Diagram]

- Cuts \( x \)-space into rectangular cells
- Works well with both categorical and quantitative features
- Interpretable result (inference)
- Decision boundary can be arbitrarily complicated

![Linear Classifier and Decision Tree Comparison]

[Comparison of linear classifiers vs. decision trees on 2 examples.]
Consider classification first. Greedy, top-down learning heuristic:

- This algorithm is more or less obvious, and has been rediscovered many times. It's naturally recursive. I'll show how it works for classification first; later I'll talk about how it works for regression.

Let $S \subseteq \{1, 2, \ldots, n\}$ be set of sample point indices.

**Top-level call:** $S = \{1, 2, \ldots, n\}$.

**GrowTree($S$)**

- If $(y_i = C$ for all $i \in S$ and some class $C$) then {
  - return new leaf($C$) [We say the leaves are pure]
- } else {
  - choose best splitting feature $j$ and splitting value $\beta$ (*)
  - $S_l = \{i \in S : X_{ij} < \beta\}$ [Or you could use $\leq$ and $>$]
  - $S_r = \{i \in S : X_{ij} \geq \beta\}$
  - return new node($j, \beta, \text{GrowTree}(S_l), \text{GrowTree}(S_r)$)

(*) How to choose best split?

- Try all splits. [All features, and all splits within a feature.]
- For a set $S$, let $J(S)$ be the cost of $S$.
- Choose the split that minimizes $J(S_l) + J(S_r)$; or, the split that minimizes weighted average $\frac{|S_l|J(S_l) + |S_r|J(S_r)}{|S|}$. [Here, I'm using the vertical brackets $|·|$ to denote set cardinality.]

How to choose cost $J(S)$?

- I'm going to start by suggesting a mediocre cost function, so you can see why it's mediocre.

**Idea 1 (bad):** Label $S$ with the class $C$ that labels the most points in $S$.

$$J(S) \leftarrow \text{# of points in } S \text{ not in class } C.$$

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- For a set $S$, let $J(S)$ be the cost of $S$.
- Choose the split that minimizes $J(S_l) + J(S_r)$; or, the split that minimizes weighted average $\frac{|S_l|J(S_l) + |S_r|J(S_r)}{|S|}$. [Here, I’m using the vertical brackets $|·|$ to denote set cardinality.]

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Problem: $J(S_l) + J(S_r) = 10$ for both splits, but left split is much better. Weighted avg prefers right split!

[There are many different splits that all have the same total cost. We want a cost function that better distinguishes between them.]
Idea 2 (good): Measure the entropy. [An idea from information theory.]
Let $Y$ be a random class variable, and suppose $P(Y = C) = p_C$. [Always nonnegative.]
The surprise of $Y$ being class $C$ is $-\log_2 p_C$. 
– event w/prob. 1 gives us zero surprise.
– event w/prob. 0 gives us infinite surprise!

[In information theory, the surprise is equal to the expected number of bits of information we need to transmit which events happened to a recipient who knows the probabilities of the events. Often this means using fractional bits, which may sound crazy, but it makes sense when you’re compiling lots of events into a single message; e.g., a sequence of biased coin flips.]

The entropy of an index set $S$ is the average surprise [when you draw a point at random from $S$],

$$H(S) = -\sum_C p_C \log_2 p_C,$$
where $p_C = \frac{|\{i \in S : y_i = C\}|}{|S|}$. [The proportion of points in $S$ that are in class C.]

If all points in $S$ belong to same class? $H(S) = -1 \log_2 1 = 0$.
Half class C, half class D? $H(S) = -0.5 \log_2 0.5 - 0.5 \log_2 0.5 = 1$.
$n$ points, all different classes? $H(S) = -\log_2 \frac{1}{n} = \log_2 n$.

[The entropy is the expected number of bits of information we need to transmit to identify the class of a sample point in $S$ chosen uniformly at random. It makes sense that it takes 1 bit to specify C or D when each class is equally likely. And it makes sense that it takes $\log_2 n$ bits to specify one of $n$ classes when each class is equally likely.]

[Left: plot of the entropy $H(p_C)$ when there are only two classes. The probability of the second class is $p_D = 1 - p_C$, so we can plot the entropy with just one dependent variable. Right: plot of the entropy $H(p_C, p_D)$ when there are three classes. The probability of the third class is $p_E = 1 - p_C - p_D$. Observe that the entropy is strictly concave.]
Weighted avg entropy after split is $H_{after} = \frac{|S_l|H(S_l) + |S_r|H(S_r)}{|S_l| + |S_r|}$.

Choose split that maximizes information gain $H(S) - H_{after}$. [Which is just the same as minimizing $H_{after}$.]
[By the way, the entropy is not the only function that works well. Many concave functions work fine, including the simple polynomial $p(1 - p)$.

More on choosing a split:
- For binary feature $x_i$: children are $x_i = 0 \& x_i = 1$.
- If $x_i$ has 3+ discrete values: split depends on application. [Sometimes it makes sense to use multiway splits; sometimes binary splits.]
- If $x_i$ is quantitative: sort $x_i$ values in $S$; try splitting between each pair of unequal consecutive values. [We can radix sort the points in linear time, and if $n$ is huge we should.]

Clever bit: As you scan sorted list from left to right, you can update entropy in $O(1)$ time per point! [This is important for obtaining a fast tree-building time.]

[Draw a row of C’s and X’s; show how we update the # of C’s and # of X’s in each of $S_l$ and $S_r$ as we scan from left to right.]

Algs & running times:
- Classify test point: Walk down tree until leaf. Return its label. Worst-case time is $O(\text{tree depth})$.
  - For binary features, that’s $\leq d$. [Quantitative features may go deeper.]
  - Usually (not always) $\leq O(\log n)$.
- Training: For binary features, try $O(d)$ splits at each node.
  - For quantitative features, try $O(n’d)$ splits; $n’ = \text{points in node}$.
  - Either way $\Rightarrow O(n’d)$ time at this node.

[Quantitative features are asymptotically just as fast as binary features because of our clever way of computing the entropy for each split.]

Each point participates in $O(\text{depth})$ nodes, costs $O(d)$ time in each node. [This is an amortized analysis: we are charging $O(d \text{ depth})$ time to each sample point.]

Running time $\leq O(\text{nd depth})$.

[As $nd$ is the size of the design matrix $X$, and the depth is often logarithmic, this is a surprisingly reasonable running time.]

---

Let $C$ be the number of class C sample points to the left of a potential split and $c$ be the number to the right of the split. Let $D$ be the number of class not-C points to the left of the split and $d$ be the number to the right of the split. Update $C$, $c$, $D$, and $d$ at each split (in $O(1)$ time per split) as you move from left to right. At each potential split, calculate the entropy of the left set as $-\frac{C}{C+D} \log_2 \frac{C}{C+D} - \frac{D}{C+D} \log_2 \frac{D}{C+D}$ and the entropy of the right set as $-\frac{c}{c+d} \log_2 \frac{c}{c+d} - \frac{d}{c+d} \log_2 \frac{d}{c+d}$. Note: log 0 is undefined, but this formula works if we use the convention $0 \log 0 = 0$.

It follows that the weighted average of the two entropies is $-\frac{1}{n’} \left( C \log_2 \frac{C}{C+D} + D \log_2 \frac{D}{C+D} + c \log_2 \frac{c}{c+d} + d \log_2 \frac{d}{c+d} \right)$, where $n’ = C + D + c + d$ is the total number of sample points stored in this treenode. Choose the split that minimizes this weighted average.
15 More Decision Trees, Ensemble Learning, and Random Forests

DECISION TREE VARIATIONS

[Last lecture, I taught you the vanilla algorithms for building decision trees and using them to classify test points. There are many variations on this basic algorithm; I'll discuss a few now.]

Decision Tree Regression

Creates a piecewise constant regression fn. [This seems too rudimentary to be true, but it’s true.]

Leaf stores label \( \mu_S = \frac{1}{|S|} \sum_{i \in S} y_i \), the mean label for training pts \( i \in S \).

Cost \( J(S) = \text{Var}((y_i : i \in S)) = \frac{1}{|S|} \sum_{i \in S} (y_i - \mu_S)^2 \).

[So if all the points in a leaf have the same \( y \)-value, then the cost is zero.]
[We choose the split that minimizes the weighted average of the variances of the children after the split.]
Stopping Early

[The basic version of the decision tree algorithm keeps subdividing treenodes until every leaf is pure. We don’t have to do that; sometimes we prefer to stop subdividing treenodes earlier.]

Why?
- Limit tree depth (for speed)
- Limit tree size (big data sets)
- Pure tree may overfit
- Given noise or overlapping distributions, pure leaves tend to overfit; better to stop early and estimate posterior probs

[When you have strongly overlapping class distributions, refining the tree down to one training point per leaf is absolutely guaranteed to overfit, giving you a classifier akin to the 1-nearest neighbor classifier. It’s better to stop early, then classify each leaf node by taking a vote of its training points; this gives you a classifier akin to a k-nearest neighbor classifier.]

[Alternatively, you can use the points to estimate a posterior probability for each leaf, and return that. If there are many points in each leaf, the posterior probabilities might be reasonably accurate.]

In the decision tree at left, each leaf has multiple classes. Instead of returning the majority class, each leaf could return a posterior probability histogram, as illustrated at right.

Leaves with multiple points return
- a majority vote or class posterior probs (classification) or
- an average (regression).

How to stop? Select stopping condition(s):
- Next split doesn’t reduce entropy/error enough (dangerous; pruning is better)
- Most of node’s points (e.g., > 90%) have same class [to deal with outliers & overlapping distribs]
- Node contains few training points (e.g., < 10) [especially for big data]
- Box’s edges are all tiny [super-fine resolution may be pointless]
- Depth too great [risky if there are still many training points in the box]
- Use validation to compare

[The last is the slowest but most effective way to know when to stop: use validation to decide whether splitting the node lowers your validation error. But if your goal is to avoid overfitting, it’s generally even more effective to grow the tree a little too large and then use validation to prune it back . . . ]
Pruning

Grow tree too large; greedily remove each split whose removal improves validation performance. [We have to do validation once for each split that we’re considering reversing.] More reliable than stopping early.

The reason why pruning often works better than stopping early is because often a split that doesn’t seem to make much progress is followed by a split that makes a lot of progress. If you stop early, you’ll never find out. Pruning is simple and highly recommended when you have enough time to build and prune the tree.

At left, a plot of decision tree leaf nodes vs. errors for baseball hitter data. At right, the best decision tree has three leaves. Players’ salaries: R1 = $165,174, R2 = $402,834, R3 = $845,346.

In this example, a 10-leaf decision tree was constructed to predict the salaries of baseball players, based on their years in the league and average hits per season. Then the tree was pruned by validation. The best decision tree on the validation data turned out to have just three leaves.

Observe that this tree is very interpretable. You could easily explain it to grandpa.

It might seem expensive to do validation once for each split we consider reversing. But you can do it pretty cheaply. What you don’t do is reclassify every training point from scratch. Instead, you first compute which leaf each validation point winds up in, then for each leaf you make a list of its validation points. When you are deciding whether to remove a split, you just look at the validation points in the two leaves you’re thinking of removing, and see how they will be reclassified and how that will change the error rate. You can compute this very quickly.

After we determine which leaf boxes each validation point ends up in, we find that pruning these two leaves improves the validation accuracy. Box colors indicate the majority classes of the training points (not shown). Local validation accuracy improves from 7/16 to 9/16.
Multivariate Splits

Find non-axis-aligned splits with other classification algs or by generating them randomly.

[An example where an ordinary decision tree needs many splits to approximate a diagonal linear decision boundary, but a single multivariate split takes care of it.]

[Here you can use other classification algorithms such as SVMs, logistic regression, and Gaussian discriminant analysis. Decision trees permit these algorithms to find nonlinear decision boundaries by making them hierarchical.]

May gain better classifier at cost of worse interpretability or speed.

[Standard decision trees are very fast because they check only one feature at each treenode. But if there are hundreds of features, and you have to check all of them at every level of the tree to classify a point, it slows down classification a lot.]

[A good compromise is to set a limit on the number of features you check at each treenode—say, three. You can use forward stepwise selection at each treenode to choose the three features.]

[On exams, assume we check only one feature per treenode unless we say otherwise!]

ENSEMBLE LEARNING

Decision trees are fast, simple, interpretable, easy to explain, invariant under scaling/translation, robust to irrelevant features.

But not the best at prediction. [Compared to previous methods we’ve seen.]

High variance. [Though we can achieve very low bias.]

[For example, suppose we take a training data set, split it into two halves, and train two decision trees, one on each half of the data. It’s not uncommon for the two trees to turn out very different. In particular, if the two trees pick different features for the very first split at the root of the tree, then it’s quite common for the trees to be completely different. So decision trees tend to have high variance.]

[So let’s think about how to fix this. As an analogy, imagine that you are generating random numbers from some distribution. If you generate just one number, it might have high variance. But if you generate $n$ random numbers and take their average, then the variance of that average is $n$ times smaller. So you might ask yourself, can we reduce the variance of decision trees by taking an average answer of a bunch of decision trees? Yes we can.]
[A 1906 county fair in Plymouth, England had a contest to guess the weight of an ox. A scientist named Francis Galton was there, and he did an experiment. He calculated the median of everyone’s guesses. The median guess was 1,207 pounds, and the true weight was 1,198 pounds, so the error was less than 1%. Even the cattle experts present didn’t estimate it that accurately.]

[National Public Radio repeated the experiment in 2015 with a cow named Penelope whose photo they published online. They got 17,000 guesses, and the average guess was 1,287 pounds. Penelope’s actual weight was 1,355 pounds, so the crowd got it to within 5 percent.]

[The main idea is that sometimes the average opinion of a bunch of idiots is better than the opinion of one expert. And so it is with learning algorithms. We call a learning algorithm a weak learner if it does better than guessing randomly. And we combine a bunch of weak learners to get a strong one.]

We can take average of output of

- different learning alg
- same learning alg on many training sets  [if we have tons of data]
- bagging: same learning alg on many random subsamples of one training set
- random forests: randomized decision trees on random subsamples

[These last two are the most common ways to use averaging, because usually we don’t have enough training data to use fresh data for every learner.]

Metalearner takes test point, feeds it into all $T$ learners, returns majority class or average output.

[Averaging is not specific to decision trees; it can work with many different learning algorithms. But it works particularly well with decision trees.]

Regression alg: take median or mean output  [of all the weak learners]
Classification alg: take majority vote OR average posterior probs

[Apology to readers: I show some videos in this lecture, which cannot be included in this report.]

[Show averageaxis.mov] [Here’s a simple classifier that takes an average of “stumps,” trees of depth 1. Observe how good the posterior probabilities look.]
[Show averageaxistree.mov] [Here’s a 4-class classifier with depth-2 trees.]
The Netflix Prize was an open competition for the best collaborative filtering algorithm to predict user ratings for films, based on previous ratings. It ran for three years and ended in 2009. The winners used an extreme ensemble method that took an average of many different learning algorithms. In fact, a couple of top teams combined into one team so they could combine their methods. They said, “Let’s average our models and split the money,” and that’s what happened.

Use learners with low bias (e.g., deep decision trees).

High variance & some overfitting are okay. Averaging reduces the variance!

Each learner may overfit, but each overfits in its own unique way.

Averaging sometimes reduces bias & increases flexibility a bit, but not reliably.

e.g., creating nonlinear decision boundary from linear classifiers.

Averaging rarely reduces bias as much as it reduces variance, so get the bias small before averaging.

Hyperparameter settings usually different than 1 learner. [Averaging reduces variance more than bias.]

Sometimes the number of trees is said to be a hyperparameter, but it’s not really true. When random forests work, usually the variance drops monotonically with the number of trees, and the main limit on the number of trees is computation time. The stopping conditions are the hyperparameters.

**Bagging = Bootstrap AGGregatING (Leo Breiman, 1994)**

Leo Breiman was a statistics professor right here at Berkeley. He did his best work after he retired in 1993. The bagging algorithm was published the following year, and then he went on to co-invent random forests as well. Unfortunately, he died in 2005.

Bagging is a randomized method for creating many different learners from the same data set. It works well with many different learning algorithms. One exception seems to be $k$-nearest neighbors; bagging mildly degrades it.

Given $n$-point training sample, generate random subsample of size $n'$ by sampling with replacement. Some points chosen multiple times; some not chosen.
If \( n' = n \), \( \sim 63.2\% \) are chosen. [On average; this fraction varies randomly.]

Build learner. Points chosen \( j \) times have greater weight:
[If a point is chosen \( j \) times, we want to treat it the same way we would treat \( j \) different points all bunched up infinitesimally close together.]

- Decision trees: \( j \)-time point has \( j \times \) weight in entropy.
- SVMs: \( j \)-time point incurs \( j \times \) penalty to violate margin.
- Regression: \( j \)-time point incurs \( j \times \) loss.

Repeat until \( T \) learners.

Random Forests

Random sampling isn’t random enough!
[With bagging, often the decision trees look very similar. Why is that?] One really strong predictor → same feature split at top of every tree.
[For example, if you’re building decision trees to identify spam, the first split might always be “viagra.” Random sampling might not change that. If the trees are very similar, then taking their average doesn’t reduce the variance much.]

Let’s reduce the correlation between different trees.

Idea: At each treenode, take random sample of \( m \) features (out of \( d \)).

Choose best split from \( m \) features.
[We’re not allowed to split on the other \( d - m \) features!]

Different random sample for each treenode.

\( m \approx \sqrt{d} \) works well for classification; \( m \approx d/3 \) for regression.
[So if you have a 100-dimensional feature space, you randomly choose 10 features and pick the one of those 10 that gives the best split. But \( m \) is a hyperparameter, and you might get better results by tuning it for your particular application. These values of \( m \) are good starting guesses.]

Smaller \( m \) → more randomness, less tree correlation, more bias
[One reason this works is if there’s a really strong predictor, only a fraction of the trees can choose that predictor as the first split. That fraction is \( m/d \). So the split tends to “decorrelate” the trees. And that means that when you take the average of the trees, your average will have less variance than a single tree.]

[You have to be careful, though, because you don’t want to dumb down the trees too much in your quest for decorrelation. Averaging works best when you have very strong learners that are also diverse. But it’s hard to create a lot of learners that are very different yet all very smart. The Netflix Prize winners did it, but it was a huge amount of work.]

Sometimes test error drops even at 100s or 1,000s of decision trees!

Disadvantages: slow; loses interpretability/inference.
[But the compensation is it’s more accurate than a single decision tree.]

[I will end by showing you examples of a very non-standard method for random forests that works magic in certain difficult circumstances.]

Idea: generate \( s \) random multivariate splits (oblique lines, quadrics); choose best split.
[You have to be a bit clever about how you generate random decision boundaries; I’m not going to discuss that. I’ll just show lots of examples.]
[Show treesidesdeep.mov] [Lots of good-enough conic random decision trees.]
[Show averageline.mov]
[Show averageconic.mov]
[Show square.mov] [Depth 2; look how good the posterior probabilities look.]
[Show squaresmall.mov] [Depth 2; see the uncertainty away from the center.]
[Show spiral2.mov] [Doesn’t look like a decision tree at all, does it?]
[Show overlapdepth14.mov] [Overlapping classes. This example overfits!]
[Show overlapdepth5.mov] [Better fit.]

500.pdf [Random forest classifiers for 4-class spiral data. Each forest takes the average of 400 trees. The top row uses trees of depth 4. The bottom row uses trees of depth 12. From left to right, we have axis-aligned splits, splits with lines with arbitrary rotations, and splits with conic sections. Each split is chosen to be the best of 500 random choices.]

randomness.pdf [Random forest classifiers for the same data. Each forest takes the average of 400 trees. In these examples, all the splits are axis-aligned. The top row uses trees of depth 4. The bottom row uses trees of depth 12. From left to right, we choose each split from 1, 5, or 50 random choices. The more choices, the less bias and the better the classifier.]
16 The Kernel Trick

KERNELS

Recall featurizing map $\Phi : \mathbb{R}^d \to \mathbb{R}^D$. $d$ input features; $D$ features after featurization ($\Phi$).
Degree-$p$ polynomials blow up to $D \in \Theta(d^p)$ features.
[When $d$ and $p$ are not small, this gets computationally intractable really fast. As I said in Lecture 4, if you have 100 features per feature vector and you want to use degree-4 polynomial decision functions, then each featurized feature vector has a length of roughly 4 million.]
Today, magically, we use those features without computing them!

Observation: In many learning algs,
– the weights can be written as a linear combo of training points, &
– we can use inner products of $\Phi(x)$’s only $\Rightarrow$ don’t need to compute $\Phi(x)$!

Suppose $w = X^\top a = \sum_{i=1}^n a_i X_i$ for some $a \in \mathbb{R}^n$.
Substitute this identity into alg. and optimize $n$ dual weights $a$ (aka dual parameters) instead of $D$ primal weights $w$.

Kernel Ridge Regression

Center $X$ and $y$ so their means are zero: $X_i \leftarrow X_i - \mu_X, \quad y_i \leftarrow y_i - \mu_y, \quad X_{i,d+1} = 1$ [don’t center the 1’s!]
This lets us replace $I'$ with $I$ in normal equations:

$$(X^\top X + \lambda I)w = X^\top y.$$ 

[To kernelize ridge regression, we need the weights to be a linear combination of the training points. Unfortunately, that only happens if we penalize the bias term $w_{d+1} = \alpha$, as these normal equations do. Fortunately, when we center $X$ and $y$, the “expected” value of the bias term is zero. The actual bias won’t usually be exactly zero, but it will often be close enough that we won’t do much harm by penalizing the bias term.]

Suppose $a$ is a solution to

$$(XX^\top + \lambda I)a = y. \quad [\text{Always has a solution if } \lambda > 0.]$$

Then $X^\top y = X^\top XX^\top a + \lambda X^\top a = (X^\top X + \lambda I)X^\top a$.
Therefore, $w = X^\top a$ is a solution to the normal equations, and $w$ is a linear combo of training points!

$a$ is a dual solution; solves the dual form of ridge regression:

Find $a$ that minimizes $\|XX^\top a - y\|^2 + \lambda ||X^\top a||^2$.

[We obtain this dual form by substituting $w = X^\top a$ into the original ridge regression cost function.]

Training: Solve $(XX^\top + \lambda I)a = y$ for $a$.
Testing: Regression fn is

$$h(z) = w^\top z = a^\top Xz = \sum_{i=1}^n a_i (X_i^\top z) \quad \Leftarrow \text{weighted sum of inner products}$$
Let \( k(x, z) = x^\top z \) be kernel fn.

[Later, we’ll replace \( x \) and \( z \) with \( \Phi(x) \) and \( \Phi(z) \), and that’s where the magic will happen.]

Let \( K = XX^\top \) be \( n \times n \) kernel matrix. Note \( K_{ij} = k(X_i, X_j) \).

\( K \) may be singular. If so, probably no solution if \( \lambda = 0 \). [Then we must choose a positive \( \lambda \). But that’s okay.]

Always singular if \( n > d + 1 \). [But don’t worry about the case \( n > d + 1 \), because you would only want to use the dual form when \( d > n \), i.e., for polynomial features. But \( K \) could still be singular when \( d > n \).]

Dual/kernel ridge regr. alg:

\[
\forall i, j, \quad K_{ij} \leftarrow k(X_i, X_j) \quad \Leftarrow O(n^2d) \text{ time}
\]

Solve \((K + \lambda I) a = y \) for \( a \) \( \Leftarrow O(n^3) \text{ time} \)

for each test pt \( z \)

\[
h(z) \leftarrow \sum_{i=1}^{n} a_i k(X_i, z) \quad \Leftarrow O(nd) \text{ time/test pt}
\]

Does not use \( X_i \) directly! Only \( k \). [This will become important soon.]

**Important:** dual ridge regression produces the same predictions as primal ridge regression (with a penalized bias term)! The difference is the running time; the dual algorithm is faster if \( d > n \), because the primal algorithm solves a \( d \times d \) linear system, whereas the dual algorithm solves an \( n \times n \) linear system.

**The Kernel Trick (aka Kernelization)**

[Here’s the magic part. We can compute a polynomial kernel without actually computing the features.]

The polynomial kernel of degree \( p \) is \( k(x, z) = (x^\top z + 1)^p \).

Theorem: \((x^\top z + 1)^p = \Phi(x)^\top \Phi(z)\) for some \( \Phi(x) \) containing every monomial in \( x \) of degree \( 0 \ldots p \).

Example for \( d = 2, p = 2 \):

\[
(x^\top z + 1)^2 = x_1^2z_1^2 + x_2^2z_2^2 + 2x_1z_1z_2 + 2x_1z_1 + 2x_2z_2 + 1
\]

\[
= \begin{bmatrix} x_1^2 & x_2^2 & \sqrt{2}x_1 & \sqrt{2}z_2 & \sqrt{2}z_1 \end{bmatrix} \begin{bmatrix} z_1^2 & z_2^2 & \sqrt{2}z_1z_2 & \sqrt{2}z_1 & \sqrt{2}z_2 \end{bmatrix}^\top
\]

[This is how we’re defining \( \Phi \).]

[Notice the factors of \( \sqrt{2} \). If you try a higher polynomial degree \( p \), you’ll see a wider variety of these constants. We have no control of the constants that appear in \( \Phi(x) \), but they don’t matter much, because the primal weights \( w \) will scale themselves to compensate. Even though we don’t directly compute the primal weights, they implicitly exist in the form \( w = X^\top a \).]

**Key win:** compute \( \Phi(x)^\top \Phi(z) \) in \( O(d) \) time instead of \( O(D) = O(d^p) \), even though \( \Phi(x) \) has length \( D \).

Kernel ridge regr. replaces \( X_i \) with \( \Phi(X_i) \); let \( k(x, z) = \Phi(x)^\top \Phi(z) \), but doesn’t compute \( \Phi(x) \) or \( \Phi(z) \); it computes \( k(x, z) = (x^\top z + 1)^p \).

**Running times for 3 ridge algns:**

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<th>train</th>
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<th>kernel</th>
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Kernel Logistic Regression

Let \( \Phi(X) \) be \( n \times D \) matrix with rows \( \Phi(X_i)^\top \).  \([\Phi(X) \text{ is the design matrix of the featurized training points.}]\)

Featurized logi. regr. with batch grad. descent:

\[
\begin{align*}
  w & \leftarrow 0 & \text{[starting point is arbitrary]} \\
  \text{repeat until convergence} & \\
  w & \leftarrow w + \epsilon (\Phi(X)^\top (y - s(\Phi(X)w))) & \text{apply } s \text{ component-wise to vector } \Phi(X)w \\
  \text{for each test pt } z & \\
  h(z) & \leftarrow s(w^\top \Phi(z))
\end{align*}
\]

Dualize with \( w = \Phi(X)^\top a \).

Then the code \( \{ a \leftarrow a + \epsilon (y - s(\Phi(X)w)) \} \) has same effect as \( \{ w \leftarrow w + \epsilon (\Phi(X)^\top (y - s(\Phi(X)w))) \} \).

Let \( K = \Phi(X)\Phi(X)^\top \). \([\text{The } n \times n \text{ kernel matrix; but we don’t compute } \Phi(X)\text{—we use the kernel trick.}]\)

Note that \( Ka = \Phi(X)\Phi(X)^\top a = \Phi(X)w. \) \([\text{And } \Phi(X)w \text{ appears in the algorithm above.}]\)

Dual/kernel logistic regression:

\[
\begin{align*}
  a & \leftarrow 0 & \text{[starting point is arbitrary]} \\
  \forall i, j, K_{ij} & \leftarrow k(X_i, X_j) & \left(\text{kernel trick}\right) \\
  \text{repeat until convergence} & \\
  a & \leftarrow a + \epsilon (y - s(Ka)) & \text{[apply } s \text{ component-wise]} \\
  \text{for each test pt } z & \\
  h(z) & \leftarrow s\left(\sum_{i=1}^n a_i k(X_i, z)\right) & \text{[kernel trick]}
\end{align*}
\]

[For classification, you can skip the logistic function \( s(\cdot) \) and just compute the sign of the summation.]

[Kernel logistic regression computes the same answer as the primal algorithm, but the running time changes.]

Important: running times depend on original dimension \( d \), not on length \( D \) of \( \Phi(\cdot)! \)

Training for \( j \) iterations:

Primal: \( O(nDj) \) time \quad Dual (no kernel trick): \( O(n^2D + n^2j) \) time \quad Kernel: \( O(n^2d + n^2j) \) time

Alternative training: stochastic gradient descent (SGD). Primal logistic SGD step is

\[
 w \leftarrow w + \epsilon (y_i - s(\Phi(X_i)^\top w)) \Phi(X_i).
\]

Dual logistic SGD maintains a vector \( q = Ka \in \mathbb{R}^n \). Note that \( q_i = (\Phi(X)w_i) = (\Phi(X_i)^\top w) \).

Let \( K_{ij} \) denote column \( i \) of \( K \).

\[
\begin{align*}
  a & \leftarrow 0; q \leftarrow 0; \forall i, j, K_{ij} \leftarrow k(X_i, X_j) & \text{[If you choose a different starting point, set } q \leftarrow Ka.\text{]} \\
  \text{repeat until convergence} & \\
  \text{choose random } i \in [1, n] & \\
  a_i & \leftarrow a_i + \epsilon (y_i - s(q_i)) & \text{[computes } q = Ka \text{ in } O(n) \text{ time, not } O(n^2) \text{ time]} \\
  q & \leftarrow q + \epsilon (y_i - s(q_i)) K_{si} & \text{[computes } q = Ka \text{ in } O(n) \text{ time, not } O(n^2) \text{ time]}
\end{align*}
\]

[SGD updates only one dual weight \( a_i \) per iteration; that’s a nice benefit of the dual formulation. We cleverly update \( q = Ka \) in linear time instead of performing a quadratic-time matrix-vector multiplication.]

Primal: \( O(Dj) \) time \quad Dual (no kernel trick): \( O(n^2D + nj) \) time \quad Kernel: \( O(n^2d + nj) \) time

Alternative testing: If \# of training points and test points both exceed \( D/d \), classifying with primal weights \( w \) may be faster. \([\text{This applies to ridge regression as well.}]\)

\[
\begin{align*}
  w & = \Phi(X)^\top a & \left(\text{once only}\right) \\
  \text{for each test pt } z & \\
  h(z) & \leftarrow s(w^\top \Phi(z)) & \text{[SGD updates only one dual weight per iteration; that’s a nice benefit of the dual formulation. We cleverly update } q = Ka \text{ in linear time instead of performing a quadratic-time matrix-vector multiplication.]} \\
  \end{align*}
\]

Primal: \( O(Dj) \) time \quad Dual (no kernel trick): \( O(n^2D + nj) \) time \quad Kernel: \( O(n^2d + nj) \) time

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The Gaussian Kernel

[Mind-blowing as the polynomial kernel is, I think our next trick is even more mind-blowing. Since we can now do fast computations in spaces with exponentially large dimensions, why don’t we go all the way and generate feature vectors in an infinite-dimensional space?]

**Gaussian kernel**, aka radial basis fn kernel: there exists a \( \Phi : \mathbb{R}^d \rightarrow \mathbb{R}^\infty \) such that

\[
k(x, z) = \exp\left(-\frac{\|x - z\|^2}{2\sigma^2}\right)
\]

[This kernel takes \( O(d) \) time to compute.]

[In case you’re curious, here’s the feature vector that gives you this kernel, for the case where you have only one input feature per sample point.]

e.g., for \( d = 1 \),

\[
\Phi(x) = \exp\left(-\frac{x^2}{2\sigma^2}\right) \left[1, \frac{x}{\sigma\sqrt{1!}}, \frac{x^2}{\sigma^2\sqrt{2!}}, \frac{x^3}{\sigma^3\sqrt{3!}}, \ldots\right]^T
\]

[This is an infinite vector, and \( \Phi(x) \cdot \Phi(z) \) is a series that converges to \( k(x, z) \). Nobody actually uses this value of \( \Phi(x) \) directly, or even cares about it; they just use the kernel function \( k(\cdot, \cdot) \).]

[At this point, it’s best *not* to think of points in a high-dimensional space. It’s no longer a useful intuition. Instead, think of the kernel \( k \) as a measure of how similar or close together two points are to each other.]

Key observation: hypothesis \( h(z) = \sum_{j=1}^n a_j k(X_j, z) \) is a linear combo of Gaussians centered at training pts.

[The dual weights are the coefficients of the linear combination.]

[The Gaussians are a basis for the hypothesis.]

[gausskernel.pdf] A hypothesis \( h \) that is a linear combination of Gaussians centered at four training points, two with positive weights and two with negative weights. If you use ridge regression with a Gaussian kernel, your “linear” regression will look something like this.]
Very popular in practice! Why?
- Gives very smooth $h$ [In fact, $h$ is infinitely differentiable; it’s $C^\infty$-continuous.]
- Behaves somewhat like $k$-nearest neighbors, but smoother
- Oscillates less than polynomials (depending on $\sigma$)
- $k(x, z)$ interpreted as a similarity measure. Maximum when $z = x$; goes to 0 as distance increases.
- Training points “vote” for value at $z$, but closer points get weightier vote.

[The “standard” kernel $k(x, z) = x \cdot z$ assigns more weight to training point vectors that point in roughly the same direction as $z$. By contrast, the Gaussian kernel assigns more weight to training points near $z$.]

Choose $\sigma$ by validation.
$\sigma$ trades off bias vs. variance:
- larger $\sigma$ → wider Gaussians & smoother $h$ → more bias & less variance

[The decision boundary (solid black) of a soft-margin SVM with a Gaussian kernel. Observe that in this example, it comes reasonably close to the Bayes optimal decision boundary (dashed purple). The dashed black curves are the boundaries of the margin. The small black disks are the support vectors that lie on the margin boundary.]

[By the way, there are many other kernels that, like the Gaussian kernel, are defined directly as kernel functions without worrying about $\Phi$. But not every function can be a kernel function. A function is qualified only if it always generates a positive semidefinite kernel matrix, for every sample. There is an elaborate theory about how to construct valid kernel functions. However, you probably won’t need it. The polynomial and Gaussian kernels are the two most popular by far.]

[As a final word, be aware that not every featurization $\Phi$ leads to a kernel function that can be computed faster than $\Theta(D)$ time. In fact, the vast majority cannot. Featurizations that can are rare and special.]
17 Neural Networks

NEURAL NETWORKS

Can do both classification & regression.

[They tie together several ideas from the course: perceptrons, linear regression, logistic regression, ensembles of learners, and stochastic gradient descent. They also tie in the idea of lifting sample points to a higher-dimensional feature space, but with a new twist: neural nets can learn features themselves.]

[I want to begin by reminding you of the story I told you at the beginning of the semester, about Frank Rosenblatt's invention of perceptrons in 1957. Remember that he held a press conference where he predicted that perceptrons would be “the embryo of an electronic computer that [the Navy] expects will be able to walk, talk, see, write, reproduce itself and be conscious of its existence.”]

[Perceptron research continued until something monumental happened in 1969. Marvin Minsky, one of the founding fathers of AI, and Seymour Papert published a book called “Perceptrons.” Sounds promising? Well, part of the book was devoted to things perceptrons can’t do. And one of those things is XOR.]

\[
\begin{array}{c|c|c}
\text{XOR} & 0 & 1 \\
\hline
0 & 0 & 1 \\
1 & 1 & 0 \\
\end{array}
\]

[Think of the four outputs here as sample points in two-dimensional space. Two of them are in class 1, and two of them are in class 0. We want to find a linear classifier that separates the 1’s from the 0’s. Can we do it? No.]

[So Minsky and Papert were basically saying, “Frank. You’re telling us this machine is going to be conscious of its own existence but it can’t do XOR?”]

[The book had a devastating effect on the field. After its publication, almost no research was done on neural net-like ideas for a decade, a time we now call the “AI Winter.” Shortly after the book was published, Frank Rosenblatt died.]

[One thing I don’t understand is why the book was so fatal when there are several almost obvious ways to get around the XOR problem. Here’s the easiest.]

If you add one new quadratic feature, \(x_1 x_2\), XOR is linearly separable in 3D.

[Draw this by hand. xorcube.pdf]

[Now we can find a plane that cuts through the cube obliquely and separates the 0’s from the 1’s.]
[However, there’s an even more powerful way to do XOR. The idea is to design linear classifiers whose output is the input to other linear classifiers. That way, you should be able to emulate arbitrary logic circuits. Suppose I put together some linear decision functions like this.]

\[
\begin{align*}
\text{linear combo} &\quad \rightarrow \quad \text{linear combo} \\
\text{linear combo} &\quad \rightarrow \quad z \\
x_1 \quad \quad &\quad \rightarrow \quad \text{linear combo} \\
x_2 \quad \quad &\quad \rightarrow \quad \text{linear combo}
\end{align*}
\]

[Draw this by hand. \[\text{lincombo.pdf}\]]

[If I interpret the output as true if \(z\) is greater than one-half or false if \(z\) is less than one-half, can I do XOR with this?]

A linear combo of a linear combo is a linear combo . . . only works for linearly separable points.

[We need one more idea to make neural nets. We need to add some sort of nonlinearity between the linear combinations. Let’s call these boxes that compute linear combinations “neurons.” If a neuron sends the linear combination it computes through some nonlinear function before sending it on to other neurons, then the neurons can act somewhat like logic gates. The nonlinearity could be as simple as clamping the output so it can’t go below zero. And that’s what people usually use in practice these days.]

[However, the traditional choice has been to use the logistic function. The logistic function can’t go below zero or above one, which is nice because it can’t ever get huge and oversaturate the other neurons it’s sending information to. The logistic function is also smooth, which means it has well-defined gradients and Hessians we can use for optimization. And we know that the logistic is often a good model for posterior probabilities.]

[With logistic functions between the linear combinations, here’s a two-level perceptron that computes the XOR function.]

\[
\begin{align*}
\text{NAND} &\quad \rightarrow \quad s(30 - 20x_1 - 20x_2) \\
\text{AND} &\quad \rightarrow \quad s(20v + 20w - 30) \\
\text{OR} &\quad \rightarrow \quad s(20x_1 + 20x_2 - 10)
\end{align*}
\]

[Draw this by hand. \[\text{xorgates.pdf}\]]

[A natural question is: can an algorithm learn a function like this?]
Network with 1 Hidden Layer

Input layer: \( x_1, \ldots, x_d; x_{d+1} = 1 \)
Hidden units: \( h_1, \ldots, h_m; h_{m+1} = 1 \)
Output layer: \( z_1, \ldots, z_k \)

Layer 1 weights: \( m \times (d + 1) \) matrix \( V \)
Layer 2 weights: \( k \times (m + 1) \) matrix \( W \)

Recall \([\text{logistic function}]\) \( s(\gamma) = \frac{1}{1 + e^{-\gamma}}\). Other nonlinear fns can be used, called the \( \text{activation fns} \).

For vector \( v \), \( s(v) = \begin{bmatrix} s(v_1) \\ s(v_2) \\ \vdots \\ 1 \end{bmatrix}, s_1(v) = \begin{bmatrix} s(v_1) \\ s(v_2) \\ \vdots \\ 1 \end{bmatrix} \) [We apply \( s \) to a vector component-wise.]

\[
\begin{align*}
  h &= s_1(Vx) \quad \text{... that is, } h_i = s(V_i^T x) \\
  z &= s(Wh) = s(Ws_1(Vx))
\end{align*}
\]

[Neural networks often have more than one output. This allows us to build multiple classifiers that share hidden units. One of the interesting advantages of neural nets is that if you train multiple classifiers simultaneously, sometimes some of them come out better because they can take advantage of particularly useful hidden units that first emerged to support one of the other classifiers.]

[We can add more hidden layers, and for image recognition tasks it’s common to have 8 to 200 hidden layers. There are many variations you can experiment with—for instance, you can have connections that go forward more than one layer.]
Training

Usually stochastic or batch gradient descent.

Pick loss fn \( L(z, y) \) e.g., \( L(z, y) = \|z - y\|^2 \)

\[ \uparrow \quad \uparrow \]
predictions true labels (could be vectors)

Cost fn is \( J(V, W) = \frac{1}{n} \sum_{i=1}^{n} L(z(X_i), Y_i) \)

[I’m using a capital \( Y \) here because now \( Y \) is a matrix with one row for each sample point and one column for each output unit of the neural net. Sometimes there is just one output unit, but many neural net applications have more.]

[Now we want to find the weight matrices \( V \) and \( W \) that minimize \( J \).]

Usually there are many local minima!

[The cost function for a neural net is, generally, not even close to convex. Sometimes, it’s possible to wind up in a bad minimum. We’ll talk later about some clever ways to coax neural nets into better minima.]

[Now let me ask you this. Suppose we start by setting all the weights to zero, and then we do gradient descent on the weights. What will go wrong?]

[This neural network has a symmetry: there’s really no difference between one hidden unit and any other hidden unit. The gradient descent algorithm has no way to break the symmetry between hidden units. You can get stuck in a situation where all the weights out of an input unit have the same value, and all the weights into an output unit have the same value, and they have no way to become different from each other. To avoid this problem, and in the hopes of finding a better local minimum, we start with random weights.]

Let \( w \) be a vector containing all the weights in \( V \) & \( W \). Batch gradient descent:

\[
\begin{align*}
    w & \leftarrow \text{vector of random weights} \\
    \text{repeat} & \\
    w & \leftarrow w - \epsilon \nabla J(w)
\end{align*}
\]

[We’ve just rewritten all the weights as a vector for notational convenience. When you actually write the code, for the sake of speed, you should probably operate directly on the weight matrices \( V \) and \( W \).]

[It’s important to make sure the random weights aren’t too big, because if a unit’s output gets too close to zero or one, it can get “stuck,” meaning that a modest change in the input values causes barely any change in the output value. Stuck units tend to stay stuck. I’ll say more about that next lecture.]

[Instead of batch gradient descent, we can use stochastic gradient descent, which means we use the gradient of one sample point’s loss function at each step. Typically, we shuffle the points in a random order, or just pick one randomly at each step.]

[The hard part of this algorithm is computing the gradient. If you simply derive one derivative for each weight, you’ll find that for a network with many layers of hidden units, it takes time linear in the number of edges in the neural network to compute a derivative for one weight. Multiply that by the number of weights. We’re going to spend the rest of this lecture learning to improve the running time to linear in the number of edges.]

Naive gradient computation: \( O(\text{edges}^2) \) time
Backpropagation: \( O(\text{edges}) \) time
Computing Gradients for Arithmetic Expressions

[Let’s see what it takes to compute the gradient of an arithmetic expression. It turns into repeated applications of the chain rule from calculus.]

\[
\begin{align*}
\frac{\partial f}{\partial a} &= \frac{\partial f}{\partial d} \frac{\partial d}{\partial a} = \frac{\partial f}{\partial d} \\
\frac{\partial f}{\partial b} &= \frac{\partial f}{\partial e} \frac{\partial e}{\partial b} = \frac{\partial f}{\partial e} \\
\frac{\partial f}{\partial c} &= \frac{\partial f}{\partial e} \frac{\partial e}{\partial c} = \frac{\partial f}{\partial e}
\end{align*}
\]

Goal: compute \( \nabla f = \begin{bmatrix} 1 \\ 2 \\ 3 \end{bmatrix} \)

\[
\begin{align*}
da &= a + b \\
d &= a + b \\
e &= cd \\
e &= cd \\
f &= e^2 \\
f &= e^2 \quad \frac{\partial f}{\partial f} = 2e
\end{align*}
\]

Each value \( z \) gives partial derivative of the form

\[
\frac{\partial f}{\partial z} = \frac{\partial f}{\partial n} \frac{\partial n}{\partial z}
\]

where \( z \) is an input to \( n \).

[Draw this by hand. \text{gradientsarith.pdf} Draw the black diagram first. Then the goal (upper right). Then the green and red expressions, from left to right, leaving out the green arrows. Then the green arrows, starting at the right side of the page and moving left. Lastly, write the text at the bottom. (Use the same procedure for the next two figures.)]
What if a unit’s output goes to more than one unit? Then we need to understand a more complicated version of the chain rule. This is a standard rule of multivariate calculus:

$$\frac{\partial}{\partial \tau} L(z_1(\tau), z_2(\tau)) = \frac{\partial L}{\partial z_1} \frac{\partial z_1}{\partial \tau} + \frac{\partial L}{\partial z_2} \frac{\partial z_2}{\partial \tau} = \nabla z L \cdot \frac{\partial z}{\partial \tau}$$

With this rule, let’s compute gradients for an expression from least-squares linear regression.

$$\frac{\partial}{\partial w_1} = \frac{\partial L}{\partial z_1} \frac{\partial z_1}{\partial w_1} + \frac{\partial L}{\partial z_2} \frac{\partial z_2}{\partial w_1} = X_{11} \frac{\partial L}{\partial z_1} + X_{12} \frac{\partial L}{\partial z_2}$$

$$\frac{\partial}{\partial w_2} = \frac{\partial L}{\partial z_1} \frac{\partial z_1}{\partial w_2} + \frac{\partial L}{\partial z_2} \frac{\partial z_2}{\partial w_2} = X_{12} \frac{\partial L}{\partial z_1} + X_{22} \frac{\partial L}{\partial z_2}$$

$$\frac{\partial}{\partial \alpha} = \frac{\partial L}{\partial z_1} \frac{\partial z_1}{\partial \alpha} + \frac{\partial L}{\partial z_2} \frac{\partial z_2}{\partial \alpha} = \frac{\partial L}{\partial z_1} + \frac{\partial L}{\partial z_2}$$

[Draw this by hand.]

Observe that we’re doing dynamic programming here. We’re computing the solutions of subproblems, then using each solution to compute the solutions of several bigger problems.

In one sense, all we’ve done here is to find a longer derivation that the gradient of the least-squares regression cost function is $\nabla_w L = 2X^\top (z - y)$, where $z = Xw$. But the way we’ve divided it into a forward pass and a backward pass gives us a way to generalize it by adding more layers of computations.
The Backpropagation Alg.

[Backpropagation is a dynamic programming algorithm for computing the gradients we need to do neural net stochastic gradient descent in time linear in the number of weights.]

$V_i^T$ is row $i$ of weight matrix $V$ [and likewise for rows of $W$]

Recall $s'(\gamma) = s(\gamma) (1 - s(\gamma))$

$h_i = s(V_i \cdot x)$, so $\nabla V_i h_i = s'(V_i \cdot x) x = h_i (1 - h_i) x$

$z_j = s(W_j \cdot h)$, so $\nabla W_j z_j = s'(W_j \cdot h) h = z_j (1 - z_j) h$

$\nabla_h z_j = z_j (1 - z_j) W_j$

[Here is the arithmetic expression for the same neural network I drew for you three illustrations ago. It looks very different when you depict it like this, but don’t be fooled; it’s exactly the same network I started with. But now we treat the weights $V$ and $W$ as the inputs, rather than the point $x$.]

$W$

$\nabla_w L = \frac{\partial L}{\partial z_j} \nabla W_j z_j$

$= \frac{\partial L}{\partial z_j} z_j (1 - z_j) h$

$V$

$\nabla_v L = \frac{\partial L}{\partial h} \nabla_V h_i$

$= \frac{\partial L}{\partial h} h_i (1 - h_i) x$

$\nabla_h L = \sum_{j=1}^k \frac{\partial L}{\partial z_j} \nabla_h z_j$

$= \sum_j z_j (1 - z_j) \frac{\partial L}{\partial z_j} W_j$

Compute $\nabla_V L$, $\nabla_W L$ one row at a time.

[Draw this by hand. backpropalg.pdf]
18 Vanishing Gradients; ReLUs; Output Units and Losses; Neurobiology

THE VANISHING GRADIENT PROBLEM; ReLUs

Problem: When unit output $s$ is close to 0 or 1 for most training points, $s' = s(1 - s) \approx 0$, so gradient descent changes $s$ very slowly. Unit is “stuck.” Slow training.

[Draw flat spots, “linear” region, & maximum curvature points (at $s(\lambda) = 0.21$ and $s(\lambda) = 0.79$) of the sigmoid function. Ideally, we would stay away from the flat spots.]

[This is called the vanishing gradient problem. The more layers your network has, the more problematic this problem becomes. Most of the early attempts to train deep, many-layered neural nets failed.]

Solution: Replace sigmoids with ReLUs: rectified linear units.

ramp fn: $r(\gamma) = \max\{0, \gamma\}$.

$$r'(\gamma) = \begin{cases} 1, & \gamma \geq 0, \\ 0, & \gamma < 0. \end{cases}$$

[The derivative is not defined at zero, but we just pretend it is.]

Most neural networks today use ReLUs for the hidden units.

[However, it is still common to use sigmoids for the output units in classification problems.]

[One nice thing about ramp functions is that they and their gradients are very fast to compute. Computers compute exponentials slowly. Even though ReLUs are linear in each half of their range, they’re still nonlinear enough to easily compute functions like XOR.]

[Of course, if you replace sigmoids with ReLUs, you have to change the derivation of backprop to reflect the changes in the gradients.]

[The gradient is sometimes zero, so you might wonder if ReLUs can get stuck too. Fortunately, it’s rare for a ReLU’s gradient to be zero for all the training data; it’s usually zero for just some training points. But yes, ReLUs sometimes get stuck too; just not as often as sigmoids.]

[The output of a ReLU can be arbitrarily large, creating the danger that it might overwhelm units downstream. This is called the “exploding gradient problem,” and it is not a big problem in shallow networks, but it becomes a big problem in deep or recurrent networks.]
Initializing Weights

[Recall that we initialize a neural network with random weight values to break the symmetry between hidden units. If we make those random values too small, breaking symmetry takes a long time. But if we make them too large, we may cause the vanishing gradient problem in sigmoid units, or the exploding gradient problem in ReLUs. Here’s a rule of thumb for initializing weights into both sigmoid and ReLU units.]

For unit with fan-in $\eta$, initialize each incoming edge to random weight with mean zero, std. dev. $\frac{1}{\sqrt{\eta}}$.

[The idea is that the sum of all the incoming weights has mean zero and standard deviation of 1. The bigger the fan-in of a sigmoid unit, the easier it is to saturate it. So we choose smaller random initial weights for units with a bigger fan-in. We do the same with ReLUs to prevent their outputs from getting larger with every successive layer.]

OUTPUT UNITS

(1) Regression: usually linear output unit(s)—omit activation fn. [Unless you’re regressing probabilities.]

Usually trained with squared-error loss.

[If you want your regression function to sometimes be negative, clearly the output unit cannot be a ReLU or a sigmoid. For that purpose, a simple linear output unit is popular. You could say there is no activation function, or you could say that the activation function is the identity function. A linear output unit with a squared-error loss is doing least-squares linear regression on learned features (hidden units).]

(2) Classification: to choose from $k \geq 2$ classes, use softmax units.

[A set of $k$ separate output units can do softmax regression on learned features (hidden units). E.g., in the MNIST digit recognition problem, we would have $k = 10$ softmax output units, one for each digit class.]

Let $y \in \mathbb{R}^k$ be vector of labels [for one training point, indicating its membership in the $k$ classes].

In the final layer, compute $t = Wh \in \mathbb{R}^k$ and apply softmax activation to obtain prediction $z \in \mathbb{R}^k$.

\[
\text{Softmax output is } z_i(t) = \frac{e^{t_i}}{\sum_{j=1}^{k} e^{t_j}}. \quad \text{Each } z_i \in (0, 1); \sum_{i=1}^{k} z_i = 1.
\]

[Interpret $z_i$ as an estimate of the posterior probability that the input belongs to class $i$.]

Strongly recommended: choose labels so $\sum_{i=1}^{k} y_i = 1$.

[Usually, for any single training point, people choose one label to be 1 and the others to be 0, which is called one-hot encoding. But one-hot encoding has a small disadvantage: each prediction $z_i$ can never be exactly 1 or 0. If you choose target labels such as 0.1 or 0.9, a neural network with lots of layers, all sufficiently wide, can “interpolate” the labels, meaning that $z = y$ for every training point and the cost function achieves its global minimum. (Unless there are co-located training points with different labels.) From here on, we will only assume that the target labels sum to 1.]

To fix vanishing gradient problem, use cross-entropy loss fn instead of squared error.

[This method prevents the vanishing gradient problem in output units, but it can’t be applied to hidden units.]

For $k$-class softmax output, cross-entropy is $L(z, y) = -\sum_{i=1}^{k} y_i \ln z_i$.

\[
\begin{align*}
\uparrow \text{true labels} \quad & \uparrow \text{prediction} \\
\} & \text{$k$-vectors}
\end{align*}
\]
Derivatives for $k$-class softmax backprop: [correct only if $\sum_{i=1}^{k} y_i = 1$ for each training point]

$$L(z, y) = -\sum_{i=1}^{k} y_i \ln z_i = -\sum_{i=1}^{k} y_i \left( t_i - \ln \sum_{j=1}^{k} e^{t_j} \right) = -\sum_{i=1}^{k} y_i t_i + \ln \sum_{j=1}^{k} e^{t_j},$$

$$t_i = W_i \cdot h, \quad \nabla_{W_i} t_i = h, \quad \nabla_{h} t_i = W_i,$$

$$\nabla_{W_i} L = -y_i h + \left( e^{t_i} \sum_{j=1}^{k} e^{t_j} \right) h = (z_i - y_i) h,$$

$$\nabla_{W} L = (z - y) h^{\top},$$

$$\nabla_{h} L = -\sum_{i=1}^{k} y_i W_i + \left( \sum_{i=1}^{k} e^{t_i} W_i \right) \left( \sum_{j=1}^{k} e^{t_j} \right) = -\sum_{i=1}^{k} (y_i - z_i) W_i = W^{\top} (z - y).$$

[Notice that both $\nabla_{W} L$ and $\nabla_{h} L$ are linear in the error $z - y$, a nice, surprising result! So the softmax unit with cross-entropy loss does not get stuck when the softmax derivatives are small. This is related to the fact that the cross-entropy loss goes to infinity as the predicted value $z_i$ approaches zero or one. The vanishing gradient of the softmax unit is compensated for by the exploding gradient of the cross-entropy loss.]

[Now I will show you how to perform backpropagation for ReLU hidden units, a $k$-class softmax output, the cross-entropy loss function, and $\ell_2$ regularization—which helps to reduce overfitting, just like in ridge regression. Observe that because of the simplifications we made by substituting softmax into cross-entropy, we don’t compute $\nabla_{z} L$ explicitly, but we still have to backpropagate the value of $z$ itself.]

(3) Two-class classification: softmax outputs with $k = 2$ are sigmoid outputs.

$$z_1 = s(t_1 - t_2); \quad z_2 = 1 - z_1 = s(t_2 - t_1).$$

Cross-entropy loss is $L(z, y) = -y_1 \ln z_1 - y_2 \ln z_2 = -y_1 \ln z_1 - (1 - y_1) \ln(1 - z_1).$

[You’ll recognize this as the logistic loss from logistic regression. Our derivation of the gradients for softmax units applies to sigmoid units too. You don’t necessarily have to provide the second, redundant sigmoid output $z_2$ to applications, but it’s useful internally as part of computing $\nabla_{W} L$ and $\nabla_{h} L$.]

[Cross-entropy losses are only for sigmoid and softmax outputs. For linear or ReLU outputs, cross-entropy makes no sense, but squared error loss makes sense.]
NEUROBIOLOGY

[The field of artificial intelligence started with some wrong premises. The early AI researchers attacked problems like chess and theorem proving, because they thought those exemplified the essence of intelligence. They didn’t pay much attention at first to problems like vision and speech understanding. Any four-year-old can do those things, and so researchers underestimated their difficulty.]

[Today, we know better. Computers can effortlessly beat four-year-olds at chess, but they still can’t play with toys well. We’ve come to realize that rule-based symbol manipulation is not the primary defining mark of intelligence. Even rats do computations that we’re hard pressed to match with our computers. We’ve also come to realize that these are different classes of problems that require very different styles of computation. Brains and computers have very different strengths and weaknesses, which reflect their different computing styles.]

[Neural networks are partly inspired by the workings of actual brains. Let’s take a look at a few things we know about biological neurons, and contrast them with both neural nets and traditional computation.]

- CPUs: largely sequential, nanosecond gates, fragile if gate fails
  superior for arithmetic, logical rules, perfect key-based memory

- Brains: very parallel, millisecond neurons, fault-tolerant

[Neurons are continually dying. You’ve probably lost a few since this lecture started. But you probably didn’t notice. And that’s interesting, because it points out that our memories are stored in our brains in a diffuse representation. There is no one neuron whose death will make you forget that $2 + 2 = 4$. Artificial neural nets often share that resilience. Brains and neural nets seem to superpose memories on top of each other, all stored together in the same weights, sort of like a hologram.]

[In the 1920’s, the psychologist Karl Lashley conducted experiments to identify where in the brain memories are stored. He trained rats to run a maze, and then made lesions in different parts of the cerebral cortex, trying to erase the memory trace. Lashley failed; his rats could still find their way through the maze, no matter where he put lesions. He concluded that memories are not stored in any one area of the brain, but are distributed throughout it. Neural networks, properly trained, can duplicate this property.]

superior for vision, speech, associative memory

[By “associative memory,” I mean noticing connections between things. One thing our brains are very good at is retrieving a pattern if we specify only a portion of the pattern.]

[It’s impressive that even though a neuron needs a few milliseconds to transmit information to the next neurons downstream, we can perform very complex tasks like interpreting a visual scene in a tenth of a second. This is possible because neurons run in parallel, but also because of their computation style.]

[Neural nets try to emulate the parallel, associative thinking style of brains, and they are the best techniques we have for many fuzzy problems, including most problems in vision and speech. Not coincidentally, neural nets are also inferior at many traditional computer tasks such as multiplying 10-digit numbers or compiling source code.]
- **Neuron**: A cell in brain/nervous system for thinking/communication
- **Action potential** or **spike**: An electrochemical impulse fired by a neuron to communicate w/other neurons
- **Axon**: The limb(s) along which the action potential propagates; “output”
  [Most axons branch out eventually, sometimes profusely near their ends.]
  [It turns out that squids have a very large axon they use for fast propulsion by expelling jets of water. The mathematics of action potentials was first characterized in these squid axons, and that work won a Nobel Prize in Physiology in 1963.]
- **Dendrite**: Smaller limb by which neuron receives info; “input”
- **Synapse**: Connection from one neuron’s axon to another’s dendrite
  [Some synapses connect axons to muscles or glands.]
- **Neurotransmitter**: Chemical released by axon terminal to stimulate dendrite
  [When an action potential reaches an axon terminal, it causes tiny containers of neurotransmitter, called vesicles, to empty their contents into the space where the axon terminal meets another neuron’s dendrite. That space is called the synaptic cleft. The neurotransmitters bind to receptors on the dendrite and influence the next neuron’s body voltage. This sounds incredibly slow, but it all happens in 1 to 5 milliseconds.]

You have about $10^{11}$ neurons, each with about $10^4$ synapses.
Analogies: [between artificial neural networks and brains]

- **Output of unit ↔ firing rate of neuron**
  
  [An action potential is “all or nothing”—all action potentials have the same shape and size. The output of a neuron is not signified by voltage like the output of a transistor. The output of a neuron is the frequency at which it fires. Some neurons can fire at nearly 1,000 times a second, which you might think of as a strong “1” output. Conversely, some types of neurons can go for minutes without firing. But some types of neurons never stop firing, and for those you might interpret a firing rate of 10 times per second as a “0”.]

- **Weight of connection ↔ synapse strength**

- **Positive weight ↔ excitatory neurotransmitter (e.g., glutamine)**

- **Negative weight ↔ inhibitory neurotransmitter (e.g., GABA, glycine) [Gamma aminobutyric acid.]**
  
  [A typical neuron is either excitatory at all its axon terminals, or inhibitory at all its terminals. It can’t switch from one to the other. Artificial neural nets have an advantage here.]

- **Linear combo of inputs ↔ summation**
  
  [A neuron fires when the sum of its inputs, integrated over time, reaches a high enough voltage. However, the neuron body voltage also decays slowly with time, so if the action potentials are coming in slowly enough, the neuron might not fire at all.]

- **Logistic/sigmoid fn ↔ firing rate saturation**
  
  [A neuron can’t fire more than 1,000 times a second, nor less than zero times a second. This limits its ability to overpower downstream neurons. We accomplish the same thing with the sigmoid function.]

- **Weight change/learning ↔ synaptic plasticity**
  
  [Donald] Hebb’s rule (1949): “Cells that fire together, wire together.”
  
  [This doesn’t mean that the cells have to fire at exactly the same time. But if one cell’s firing tends to make another cell fire more often, their excitatory synaptic connection tends to grow stronger. There’s a reverse rule for inhibitory connections. And there are ways for neurons that aren’t even connected to grow connections.]
  
  [There are simple computer learning algorithms based on Hebb’s rule. They can work, but they’re generally not nearly as fast or effective as backpropagation.]

[Backpropagation is one part of artificial neural networks for which the analogy is doubtful. There have been some proposals that the brain might do something vaguely like backpropagation, but it seems tenuous.]

(The following items are all spoken, not written. See the figure on the next page.)

The brain is very modular.

- The part of our brain we think of as most characteristically human is the cerebral cortex, the seat of self-awareness, language, and abstract thinking.

But the brain has a lot of other parts that take the load off the cortex.

- Our brain stem regulates functions like heartbeat, breathing, and sleep.

- Our cerebellum governs fine coordination of motor skills. When we talk about “muscle memory,” much of that is in the cerebellum, and it saves us from having to consciously think about how to walk or talk or brush our teeth, so the cortex can focus on where to walk and what to say.
Our limbic system is the center of emotion and motivation, and as such, it makes a lot of the big decisions. I sometimes think that 90% of the job of our cerebral cortex is to rationalize decisions that have already been made by the limbic system. 

Our visual cortex (in the occipital lobe) performs a lot of processing on the input from your eyes to change it into a more useful form. Neuroscientists and computer scientists are particularly interested in the visual cortex for several reasons. Vision is an important problem for computers. The visual cortex is one of the easier parts of the brain to study in lab animals. The visual cortex is largely a feedforward network with few neurons going backward, so it’s easier for us to train computers to behave like the visual cortex.

Although the brain has lots of specialized modules, one thing that’s interesting about the frontal lobe is that it seems to be made of general-purpose neural tissue that looks more or less the same everywhere, at least before it’s trained. If you experience damage to part of the frontal lobe early enough in life, while your brain is still growing, the functions will just relocate to healthy parts of the frontal lobe, and you’ll probably never notice the difference.

As computer scientists, our primary motivation for studying neurology is to try to get clues about how we can get computers to do tasks that humans are good at. But neurologists and psychologists have also been part of the study of neural nets from the very beginning. Their motivations are scientific: they’re curious how humans think, and how we can do the things we do.
19 Better Neural Network Training; Convolutional Neural Networks

HEURISTICS FOR FASTER TRAINING

[A big disadvantage of neural nets is that they take a long, long time to train compared to other classification methods we’ve studied. Here are some ways to speed them up. Unfortunately, you usually have to experiment with techniques and hyperparameters to find which ones will help with your particular application. I suggest you implement vanilla backpropagation first, usually in combination with stochastic gradient descent, and experiment with fancy heuristics only after you get that working.]

(1) Fix vanishing gradient problem. [As described in the previous lecture.]

(2) Stochastic gradient descent (SGD): faster than batch on large, redundant data sets. [Whereas batch gradient descent walks downhill on one cost function, stochastic descent takes a very short step downhill on one point’s loss function and then another short step on another point’s loss function. The cost function is the sum of the loss functions over all the sample points, so one batch step behaves similarly to \( n \) stochastic steps and takes roughly the same amount of time. But if you have many different examples of the digit “9”, they contain much redundant information, and stochastic gradient descent learns the redundant information more quickly—often much more quickly. Conversely, if the data set is so small that it encodes little redundant information, batch gradient descent is typically faster.]

One epoch presents every training point once. Training usually takes many epochs, but if sample is huge [and carries lots of redundant information], SGD can take less than one epoch.
(3) Normalizing the training pts.
   - Center each feature so mean is zero.
   - Then scale each feature so variance $\approx 1$.

Neural networks are an example of an optimization algorithm whose cost function tends to have better-conditioned Hessians if the input features are normalized, so it may converge to a local minimum faster.

[Recall these illustrations from Lecture 5. Gradient descent on a function with an ill-conditioned Hessian matrix can be slow because a large step size diverges in one direction (left) while a small step size converges slowly in another direction (center). Normalizing the data might improve the conditioning of the Hessian (right), thus speeding up gradient descent. Moreover, if you use $\ell_2$-regularization, normalization makes it penalize the features more equally.]

[Remember that the power of neural networks comes from the nonlinearity of the activation function, and the nonlinearity of a sigmoid or ReLU unit falls where the linear combination of values coming in is close to zero. Centering makes it easier for the first layer of hidden units to be in the nonlinear operating region.]

[Remember that whatever linear transformation you apply to the training points, you must later apply the same linear transformation to the test points you want to classify!]
(4) Use a small random subsample of training data to choose learning rate $\epsilon$ [or adaptive learning schedule]. Practitioners have found that the size of the training set has only a weak effect on the best choice of $\epsilon$. So use a subsample to quickly estimate a good learning rate, then apply it to all your training points.

(5) Different learning rate for each layer of weights. [Commonly we find that there are large variations in the magnitudes of the gradients in different layers of edges. A learning rate that’s just right for one layer might be much too small to make progress in another layer. It looks like the illustrations of ill-conditioning on the previous page, but each axis represents a different layer of the network.]

(6) Emphasizing schemes. [Neural networks learn the most redundant examples quickly, and the most rare examples slowly. This motivates emphasizing schemes, which repeat the rare examples more often.]
   - Stochastic: present examples from rare classes more often, or w/bigger $\epsilon$.
   - Can do the same for high-loss examples. [E.g., perceptron alg. presents only misclassified examples.]
   - Batch: cost fn is a weighted average of training pts. [Examples from rare classes can be more heavily weighted.]

[Be forewarned that emphasizing schemes can backfire if you have really bad outliers.]

(7) Acceleration schemes: Adam, AdaGrad, AMSGrad, RMSprop, momentum. [These variations of stochastic gradient descent are quite popular. They speed up stochastic gradient descent by using adaptive learning rates and descent directions that aren’t the direction of steepest descent. Look them up online if you’re curious.]

DOUBLE DESCENT

[Early neural network researchers sometimes struggled with their networks falling into bad local minima and failing to achieve low training errors. But with experience and greater computational power, we’ve discovered that these problems can usually be solved simply by adding more units to every hidden layer. We call this “making the network wider.” Sometimes you also have to add more layers. But if your layers are wide enough, and there are enough of them, a well-designed neural network can typically output exactly the correct label for every training point, which implies that you’re at a global minimum of the cost function.]

Hidden layers are wide enough + numerous enough $\Rightarrow$ network output can interpolate the label ($z = y$) for every training pt $\Rightarrow$ find global minimum of cost fn.

[Last lecture, I pointed out that if you use sigmoid or softmax output units, you can’t set the labels to exactly 1 or 0 and achieve interpolation, as sigmoid and softmax outputs are strictly between 0 and 1; but with labels like 0.1 and 0.9, interpolating the labels is a realistic goal! And the linear output units used for regression can interpolate arbitrary numbers. If you fall into a bad local minimum, your network is too small.]

[One reason it took so long to make this discovery is that researchers believed that having too many weights in a neural network would cause overfitting. It turns out that’s only half true. Empirically, we observe a phenomenon called “double descent,” illustrated below.]

[Consider the solid blue curve, showing the test error as the width of a network increases. The horizontal axis is the number of units per hidden layer. As that number increases, at left the test error exhibits the classic U-shaped bias-variance tradeoff. But when we pass the point where the network is interpolating the labels and continue to add more weights, we typically see a second “descent,” where the test error starts to decrease again and ultimately gets even lower than before! The peak in the middle of the curve tends to be larger when there is more noise in the labels. Observe that the test error continues to fall even after the training error is zero. The takeaway is, “bigger models are better.”]
HEURISTICS FOR BETTER GENERALIZATION

[Classic methods for preventing overfitting, such as $\ell_2$ regularization and ensembles of learners, sometimes help neural networks to generalize better to points they haven’t been trained on.]

1. $\ell_2$ regularization, aka weight decay.

Add $\lambda \|w\|^2$ to the cost/loss fn, where $w$ is vector of all weights in network.

[w includes all the weights in all the weight matrices, rewritten as a vector.]

[We regularize for the same reason we do it in ridge regression: we suspect that overly large weights are spurious.]

[With a neural network, it’s not clear whether penalizing the bias terms is bad or good. Penalizing the bias terms has the effect of drawing each ReLU or sigmoid unit closer to the center of its nonlinear operating region. I would suggest to try both ways and use validation to decide whether you should penalize the bias terms or not. Also, you could try using a different hyperparameter for the bias terms than the $\lambda$ you use for the other weights.]

Effect: step $\Delta w_i = -\epsilon \frac{\partial J}{\partial w_i}$ has extra term $-2\epsilon \lambda w_i$

Weight $w_i$ decays by factor $1 - 2\epsilon \lambda$ if not reinforced by training.
Write “10 hidden units + softmax + cross-entropy loss.” [Examples of 2D classification without (left) and with (right) weight decay. Observe that in the second example, the decision boundary (black) better approximates the Bayes optimal boundary (dashed purple curve).]

(2) Ensemble of neural nets. Bagging + random initial weights. 
[We saw how well ensemble learning works for decision trees. It works well for neural nets too. The combination of bagging and random initial weights helps ensure that each neural net comes out differently. Obviously, ensembles of neural nets are slow.]

(3) Dropout emulates an ensemble in one network.

[During training, we temporarily disable a random subset of the units, along with all the edges in and out of those units. It seems to work well to disable each hidden unit with probability 0.5, and to disable input units with a smaller probability. We do stochastic gradient descent and we frequently change which random subset of units is disabled. The authors claim that their method gives better generalization than $\ell_2$ regularization. It gives some of the advantages of an ensemble, but it’s faster to train. I’m leaving out many important details; see Hinton et al., “Improving neural networks by preventing co-adaptation of feature detectors.”]

[Recall Karl Lashley’s rat experiments, where he tried to make rats forget how to run a maze by introducing lesions in their cerebral cortices, and it didn’t work. He concluded that the knowledge is distributed throughout their brains, not localized in one place. Dropout is a way to strengthen this effect in neural networks.]
**CONVOLUTIONAL NEURAL NETWORKS** (ConvNets; CNNs)

[Convolutional neural nets have driven a big resurgence of interest in neural nets in the last decade. Often you’ll hear the buzzword deep learning, which refers to neural nets with many layers. All the best image recognition networks are deep and convolutional. In 2018, the Association for Computing Machinery gave the Alan M. Turing Award to Geoff Hinton, Yann LeCun, and Yoshua Bengio for their work on deep neural networks.]

Vision: inputs are images. 200 × 200 image = 40,000 pixels. If we connect them all to 40,000 hidden units → 1.6 billion connections.

Vision nets can be overparametrized: too many weights.

[As we saw in our discussion of double descent, this can give you excellent generalization to new examples. The problem with having billions of weights is that the network becomes very slow to train or even to use.]

[Remember that early in the semester, I told you that you can get better performance on the handwriting recognition task by using edge detectors. Edge detectors have two interesting properties. First, each edge detector looks at just one small part of the image. Second, the edge detection computation is the same no matter which part of the image you apply it to. Let’s apply these two properties to neural net design.]

ConvNet ideas:

(1) Local connectivity: A hidden unit (in early layer) connects only to a small patch of units in previous layer.

[This speeds up both training and classification considerably.]

(2) Shared weights: Groups of hidden units share same set of weights, called a mask aka filter aka kernel.

[No relationship to the kernels of Lecture 16.]

We learn several masks.

[Each mask operates on every patch of image.]

Masks × patches = hidden units in first hidden layer.

If one mask learns to detect edges, every patch has an edge detector.

[Because the mask that detects edges is applied to every patch.]

ConvNets exploit repeated structure in images, audio.

Convolution: the same linear transformation applied to different patches of the input by shifting.

[Shared weights are a kind of regularization, because shared weights means fewer weights. It’s unlikely that a weight will become spuriously large if it’s used in many places.]

[But shared weights have another big advantage. Suppose that gradient descent starts to develop an edge detector. That edge detector is being trained on every part of every image, not just on one spot. And that’s good, because edges appear at different locations in different images. The location no longer matters; the edge detector can learn from edges wherever they are.]

[In a neural net, you can think of hidden units as features that we learn, as opposed to features that you code up yourself. Convolutional neural nets take them to the next level by learning features from multiple patches simultaneously and then applying those features everywhere, not just in the patches where they were originally learned.]

[By the way, local connectivity was inspired by the human visual system, as well as by techniques used in image processing. However, shared weights don’t happen in biology.]
Neurologists can stick needles into individual neurons in animal brains. After a few hours the neuron dies, but until then they can record its action potentials. In this way, biologists quickly learned how some of the neurons in the retina, called retinal ganglion cells, respond to light. They have interesting receptive fields, illustrated in the slides, which show that each ganglion cell receives excitatory stimulation from receptors in a small patch of the retina but inhibitory stimulation from other receptors around it.

The signals from these cells propagate to the V1 visual cortex in the occipital lobe at the back of your skull. The V1 cells proved harder to understand. David Hubel and Torsten Wiesel of the Johns Hopkins University put probes into the V1 visual cortex of cats, but they had a very hard time getting any neurons to fire there. However, a lucky accident unlocked the secret and ultimately won them the 1981 Nobel Prize in Physiology.

The glass slide happened to be at the particular orientation the neuron was sensitive to. The neuron doesn’t respond to other orientations; just that one. So they were pretty lucky to catch that.

The simple cells act as line detectors and/or edge detectors by taking a linear combination of inputs from retinal ganglion cells.

The complex cells act as location-independent line detectors by taking inputs from many simple cells, which are location dependent.

Later researchers showed that local connectivity runs through the V1 cortex by projecting certain images onto the retina and using radioactive tracers in the cortex to mark which neurons had been firing. Those images show that the neural mapping from the retina to V1 is retinatopic, i.e., locality preserving. This is a big part of the inspiration for convolutional neural networks!

Unfortunately, as we go deeper into the visual system, layers V2 and V3 and so on, we know less and less about what processing the visual cortex does.
[ConvNets were first popularized by the success of Yann LeCun’s “LeNet 5” handwritten digit recognition software. LeNet 5 has six hidden layers! Hidden layers 1 and 3 are convolutional layers in which groups of units share weights. Layers 2 and 4 are pooling layers that make the image smaller. These are just hardcoded max-functions with no weights and nothing to train. Layers 5 and 6 are just regular layers of hidden units with no shared weights. A great deal of experimentation went into figuring out the number of layers and their sizes. At its peak, LeNet 5 was responsible for reading the zip codes on 10% of US Mail. Another Yann LeCun system was deployed in ATMs and check reading machines and was reading 10 to 20% of all the checks in the US by the late 90’s. LeCun is one of the Turing Award winners I told you about earlier.]

[Show Yann LeCun’s video LeNet5.mov, illustrating LeNet 5.]

[When ConvNets were first applied to image analysis, researchers found that some of the learned masks are edge detectors or line detectors, similar to the ones that Hubel and Wiesel discovered! This created a lot of excitement in both the artificial intelligence community and the neuroscience community. The fact that a neural net can naturally learn the same features as the mammalian visual cortex is impressive.]

[I told you two lectures ago that neural net research was popular in the 60’s, but the 1969 book Perceptrons killed interest in them throughout the 70’s. They came back in the 80’s, but interest was partly killed off a second time in the 00’s by . . . guess what? By support vector machines. SVMs work well for a lot of tasks, they’re much faster to train, and they more or less have only one hyperparameter, whereas neural nets take a lot of work to tune.]

[Neural nets are now in their third wave of popularity. The single biggest factor in bringing them back is probably big data. Thanks to the internet, we now have absolutely huge collections of images to train neural nets with, and researchers have discovered that neural nets often give better performance than competing algorithms when you have huge amounts of data to train them with. In particular, convolutional neural nets are now learning better features than hand-tuned features. That’s a recent change.]

[The event that brought attention back to neural nets was the ImageNet Image Classification Challenge in 2012. The winner of that competition was a neural net, and it won by a huge margin, about 10%. It’s called AlexNet, and it’s surprisingly similarly to LeNet 5, in terms of how its layers are structured. However, there are some new innovations that led to their prize-winning performance, in addition to the fact that the training set had 1.4 million images: they used ReLUs, dropout, and GPUs for training.]

![Architecture of AlexNet.](alexnet.pdf)
20 Unsupervised Learning and Principal Components Analysis

UNSUPERVISED LEARNING

We have sample points, but no labels!
No classes, no y-values, nothing to predict.
Goal: Discover structure in the data.

Examples:
- Clustering: partition data into groups of similar/nearby points.
- Dimensionality reduction: data often lies near a low-dimensional subspace (or manifold) in feature space; matrices have low-rank approximations.
  [Whereas clustering is about grouping similar sample points, dimensionality reduction is more about identifying a continuous variation from sample point to sample point.]
- Density estimation: fit a continuous distribution to discrete data.
  [When we use maximum likelihood estimation to fit Gaussians to sample points, that’s density estimation, but we can also fit functions more complicated than Gaussians.]

PRINCIPAL COMPONENTS ANALYSIS (PCA) (Karl Pearson, 1901)

Goal: Given sample points in $\mathbb{R}^d$, find $k$ directions that capture most of the variation. (Dimensionality reduction.)
Let $X$ be an $n \times d$ design matrix. [No fictitious dimension.]

From now on, assume $X$ is centered: $\sum_i X_i = 0$.
[We center the data in the usual way: by computing the sample mean, then subtracting the sample mean from each sample point.]

Let $w$ be a unit vector.

The orthogonal projection of point $x$ onto vector $w$ is $\tilde{x} = (x \cdot w) w$

If $w$ not unit, $\tilde{x} = \frac{x \cdot w}{\|w\|^2} w$

[The idea is that we're going to pick the best direction $w$, then project all the data down onto $w$ so we can analyze it in a one-dimensional space. Of course, we lose a lot of information when we project down from $d$ dimensions to just one. So, suppose we pick several directions. Those directions span a subspace, and we want to project points orthogonally onto the subspace. This is easy if the directions are orthogonal to each other.]
Given orthonormal directions $v_1, \ldots, v_k$, $\tilde{x} = \sum_{i=1}^k (x \cdot v_i) v_i$.

[The word “orthonormal” means they’re all mutually orthogonal and all have length 1.]

Often we want just the $k$ principal coordinates $x \cdot v_i$ in principal component space.

[Often we don’t actually want the projected point in $\mathbb{R}^d$. Sometimes we do, but often we just want to principal coordinates.]

$X^\top X$ is square, symmetric, positive semidefinite, $d \times d$ matrix. [As it’s symmetric, its eigenvalues are real.]

Let $0 \leq \lambda_1 \leq \lambda_2 \leq \ldots \leq \lambda_d$ be its eigenvalues. [sorted]

Let $v_1, v_2, \ldots, v_d$ be corresponding orthogonal unit eigenvectors. These are the principal components.

[. . . and the most important principal components will be the ones with the greatest eigenvalues. I will show you this in three different ways.]

PCA derivation 1: Fit a Gaussian to data with maximum likelihood estimation.

Choose $k$ Gaussian axes of greatest variance.

Recall that MLE estimates a covariance matrix $\hat{\Sigma} = \frac{1}{n} X^\top X$. [Presuming $X$ is centered.]

PCA Alg:

− Center $X$.

− Optional: Normalize $X$. Units of measurement different?

  − Yes: Normalize.

    [Bad for principal components to depend on arbitrary choice of scaling.]

  − No: Usually don’t.

    [If several features have the same unit of measurement, but some of them have smaller variance than others, that difference is usually meaningful.]

− Compute unit eigenvectors/values of $X^\top X$. 

\[gaussfitpca.png\] [A Gaussian fitted to sample points.]
– Choose \( k \). (Optional: based on the eigenvalue sizes.)
– For the best \( k \)-dimensional subspace, pick eigenvectors \( v_{d-k+1}, \ldots, v_d \).
– Compute the \( k \) principal coordinates \( x \cdot v_i \) of each training/test point.

[When we do this projection, we have two choices: we can project the original, un-centered training data OR we can project the centered training data. But if we do the latter, we have to translate the test data by the same vector we used to translate the training data when we centered it.]

[End of algorithm.]

\[
\text{% of variability} = \frac{\sum_{i=d-k+1}^{d} \lambda_i}{\sum_{i=1}^{d} \lambda_i}
\]

[Plot of # of eigenvectors vs. percentage of sample variance captured for a 17D data set. In this example, just 3 eigenvectors capture 70% of the variance.]

[Projection of 4D data onto a 2D subspace. Each point represents one metropolitan area. Normalized data at left; unnormalized data at right. The arrows show the four original axes projected on the two principal components. When the data are not normalized, rare occurrences like murder have little influence on the principal directions. Which is better? It depends on whether you think that low-frequency events like murder and rape should have a larger influence.]

[If you are using PCA as a preprocess for a supervised learning algorithm, there’s a more effective way to choose \( k \): validation.]
PCA derivation 2: Find direction $w$ that maximizes sample variance of projected data. [In other words, when we project the data down, we don’t want it all to bunch up; we want to keep it as spread out as possible.]

Points projected on a line. We wish to choose the orientation of the green line to maximize the sample variance of the blue points.

Find $w$ that maximizes $\text{Var}(\{\tilde{X}_1, \tilde{X}_2, \ldots, \tilde{X}_n\}) = \frac{1}{n} \sum_{i=1}^{n} \left( X_i \cdot \frac{w}{\|w\|} \right)^2 = \frac{1}{n} \frac{\|Xw\|^2}{\|w\|^2} = \frac{1}{n} \frac{w^T X^T X w}{w^T w}$

Rayleigh quotient of $X^T X$ and $w$

This fraction is a well-known construction called the Rayleigh quotient. When you see it, you should smell eigenvectors nearby. How do we maximize this?

If $w$ is an eigenvector $v_i$ of $X^T X$, Ray. quo. $= \lambda_i$

→ of all eigenvectors, $v_d$ achieves maximum variance $\lambda_d/n$.

One can show $v_d$ beats every other vector too.

Because every vector $w$ is a linear combination of eigenvectors, and so its Rayleigh quotient will be a convex combination of eigenvalues. It’s easy to prove this, but I don’t have the time. For the proof, look up “Rayleigh quotient” in Wikipedia.

So the top eigenvector gives us the best direction. But we typically want $k$ directions. After we’ve picked one direction, then we have to pick a second direction that’s orthogonal to the best direction. But subject to that constraint, we again pick the direction that maximizes the sample variance.

What if we constrain $w$ to be orthogonal to $v_d$? Then $v_{d-1}$ is optimal.

[And if we need a third direction orthogonal to $v_d$ and $v_{d-1}$, the optimal choice is $v_{d-2}$. And so on.]
PCA derivation 3: Find direction $w$ that minimizes mean squared projection distance.

[This is an animated GIF; unfortunately, the animation doesn’t work in the PDF lecture notes. Find the direction of the black line for which the sum of squares of the lengths of the red lines is smallest.]

[You can think of this as a sort of least-squares linear regression, with one subtle but important change. Instead of measuring the error in a fixed vertical direction, we’re measuring the error in a direction orthogonal to the principal component direction we choose.]

![projsq.png, projpca.png](image)

[Least-squares linear regression vs. PCA. In linear regression, the projection direction is always vertical; whereas in PCA, the projection direction is orthogonal to the projection hyperplane. In both methods, however, we minimize the sum of the squares of the projection distances.]

Find $w$ that minimizes

$$
\sum_{i=1}^{n} \|X_i - \bar{X}_i\|^2 = \sum_{i=1}^{n} \left\| X_i - \frac{X_i \cdot w}{\|w\|^2} \right\|^2 = \sum_{i=1}^{n} \left( \|X_i\|^2 - \left( X_i \cdot \frac{w}{\|w\|} \right)^2 \right) = \text{constant} - n \text{(variance from derivation 2)}.
$$

Minimizing mean squared projection distance = maximizing variance.

[From this point, we carry on with the same reasoning as derivation 2.]
Illustration of the first two principal components of the single nucleotide polymorphism (SNP) matrix for the genes of various Europeans. The design matrix has 2,541 people from these locations in Europe (right), and 309,790 SNPs per person. Each SNP is binary, so think of it as 309,790 dimensions of zero or one. The output (left) shows spots on the first two principal components where there was a high density of projected people from a particular national type. What’s amazing about this is how closely the projected genotypes resemble the geography of Europe.

**Eigenfaces**

$X$ contains $n$ images of faces, $d$ pixels each.

[If we have a 200 × 200 image of a face, we represent it as a vector of length 40,000, the same way we represent the MNIST digit data.]

Face recognition: Given a query face, compare it to all training faces; find nearest neighbor in $\mathbb{R}^d$.

[This works best if you have several training photos of each person you want to recognize, with different lighting and different facial expressions.]

Problem: Each query takes $\Theta(nd)$ time.

Solution: Run PCA on faces. Reduce to much smaller dimension $d'$. Now nearest neighbors takes $O(nd')$ time.

[ Possibly even less. We’ll talk about speeding up nearest-neighbor search at the end of the semester. If the dimension is small enough, you can sometimes do better than linear time.]

[If you have 500 stored faces with 40,000 pixels each, and you reduce them to 40 principal components, then each query face requires you to read 20,000 stored principal coordinates instead of 20 million pixels.]
Unsupervised Learning and Principal Components Analysis

Facerecaverage.jpg, facereceigen0.jpg, facereceigen119.jpg, facereceigen.jpg [Images of the eigenfaces. The “average face” is the mean used to center the data.]

eigenfaceproject.pdf [Images of a face (left) projected onto the first 4 and 50 eigenvectors, with the average face added back. These last image is blurry but good enough for face recognition.]
For best results, equalize the intensity distributions first.

![facerecequalize.jpg](Image equalization.)

[Eigenfaces are not perfect. They encode both face shape and lighting. Ideally, we would have some way to factor out lighting and analyze face shape only, but that’s harder. Some people say that the first 3 eigenfaces are usually all about lighting, and you sometimes get better facial recognition by dropping the first 3 eigenfaces.]

[Optional: Show Blanz–Vetter face morphing video (morphmod.mpg).]

[Blanz and Vetter use PCA in a more sophisticated way for 3D face modeling. They take 3D scans of people’s faces and find correspondences between peoples’ faces and an idealized model. For instance, they identify the tip of your nose, the corners of your mouth, and other facial features, which is something the original eigenface work did not do. Instead of feeding an array of pixels into PCA, they feed the 3D locations of various points on your face into PCA. This works more reliably.]
21 The Singular Value Decomposition; Clustering

The Singular Value Decomposition (SVD) [and its Application to PCA]

Problems: Computing $X^\top X$ takes $\Theta(nd^2)$ time.
$X^\top X$ is poorly conditioned $\rightarrow$ numerically inaccurate eigenvectors.
[The SVD improves both these problems.]

[Earlier this semester, we talked about the eigendecomposition of a square, symmetric matrix. Unfortunately, nonsymmetric matrices don’t eigendecompose nearly as nicely, and non-square matrices don’t have eigenvectors at all. Happily, there is a similar decomposition that works for all matrices, even if they’re not symmetric and not square.]

Fact: Every $X$ has a singular value decomposition $X = UDV^\top$. If $n \geq d$, it has the form

$$\begin{align*}
X &= U \\
D &= \begin{pmatrix}
\delta_1 & & \\
& \delta_2 & \\
& & \ddots \\
& & & \delta_d \\
\end{pmatrix} \\
V^\top &= \sum_{i=1}^d \delta_i u_i v_i^\top \\
U^\top U &= I \\
V^\top V &= I
\end{align*}$$

Diagonal entries $\delta_1, \ldots, \delta_d$ of $D$ are nonnegative singular values of $X$.

[Some of the singular values might be zero. By convention, they are never negative. The number of nonzero singular values is equal to the rank of $X$. If $X$ is a centered design matrix for sample points that all lie on a line, there is only one nonzero singular value. If the centered sample points span a subspace of dimension $r$, there are $r$ nonzero singular values and rank $X = r$.]

[If $n < d$, an SVD still exists, but now $U$ is square and $V$ is not.]

Fact: $v_i$ is an eigenvector of $X^\top X$ w/eigenvalue $\delta_i^2$.
Proof: $X^\top X = VD^2V^\top = VD^2V^\top$ which is an eigendecomposition of $X^\top X$.

[The columns of $V$ are the eigenvectors of $X^\top X$, which are the principal components we need for PCA. The SVD also tells us their eigenvalues, which are the squares of the singular values of $X$. By the way, that’s related to why the SVD is more numerically stable: the ratios between singular values are smaller than the ratios between eigenvalues.]

[If $n < d$, $V$ will omit some of the eigenvectors that have eigenvalue zero, but those are useless for PCA.]
Fact: We can find the $k$ greatest singular values & corresponding vectors in $O(ndk)$ time.
[So we can save time by computing some of the singular vectors without computing all of them.]
[There are approximate, randomized algorithms that are even faster, producing an approximate SVD in $O(nd \log k)$ time. These are starting to become popular in algorithms for very big data.]
[https://code.google.com/archive/p/redsvd/]

Important: Row $i$ of $UD$ gives the principal coordinates of sample point $X_i$ (i.e., $\forall i, \forall j, X_i \cdot v_j = \delta_{ij}$).
[So we don’t need to explicitly compute the inner products $X_i \cdot v_j$; the SVD has already done it for us.]
[Proof: $XV = UD^{-1}V = UD$, so $(XV)_{ij} = (UD)_{ij}$.]

**CLUSTERING**

Partition data into clusters so points in a cluster are more similar than across clusters.

Why?
- Discovery: Find songs similar to songs you like; determine market segments
- Hierarchy: Find good taxonomy of species from genes
- Quantization: Compress a data set by reducing choices
- Graph partitioning: Image segmentation; find groups in social networks

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<td>Red</td>
<td>Green</td>
<td>Blue</td>
<td>Light Blue</td>
</tr>
</tbody>
</table>

[4-Seam Fastball, 2-Seam Fastball, Changeup, Slider, Curveball]

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**Barry Zito**

---

[zito.pdf (from a talk by Michael Pane) $k$-means clusters that classify Barry Zito’s baseball pitches. Here we discover that there really are distinct classes of baseball pitches.]
**k-Means Clustering aka Lloyd’s Algorithm (Stuart Lloyd, 1957)**

**Goal:** Partition $n$ points into $k$ disjoint clusters.

Assign each input point $X_i$ a cluster label $y_i \in [1,k]$.

Cluster $i$’s mean is $\mu_i = \frac{1}{n_i} \sum_{j:y_j=i} X_j$, given $n_i$ points in cluster $i$.

Find $y$ that minimizes $\sum_{i=1}^{k} \sum_{j:y_j=i} \|X_j - \mu_i\|^2$ [Sum of the squared distances from points to their cluster means.]

NP-hard. Solvable in $O(nk^n)$ time. [Try every partition.]

**k-means heuristic:** Alternate between

1. $y_j$’s are fixed; update $\mu_i$’s
2. $\mu_i$’s are fixed; update $y_j$’s

Halt when step (2) changes no assignments.

[So, we have an assignment of points to clusters. We compute the cluster means. Then we reconsider the assignment. A point might change clusters if some other’s cluster’s mean is closer than its own cluster’s mean. Then repeat.]

**Step (1):** One can show (calculus) the optimal $\mu_i$ is the mean of the points in cluster $i$.

[This is easy calculus, so I leave it as a short exercise.]

**Step (2):** The optimal $y$ assigns each point $X_j$ to the closest center $\mu_i$.

[If there’s a tie, and one of the choices is for $X_j$ to stay in the same cluster as the previous iteration, always take that choice.]

[... so both steps minimize the cost function, but they don’t optimize all the variables at once.]

[2means.png] [An example of 2-means. Odd-numbered steps reassign the data points. Even-numbered steps compute new means.]
Both steps decrease objective fn unless they change nothing. 
[Therefore, the algorithm never returns to a previous assignment.] 
Hence alg. must terminate. [As there are only finitely many assignments.] 
[This argument says that Lloyd’s algorithm never loops forever. But it doesn’t say anything optimistic about the running time, because we might see $O(k^n)$ different assignments before we halt. In theory, one can actually construct point sets in the plane that take an exponential number of iterations, but those don’t come up in practice.] 
Usually very fast in practice. Finds a local minimum, often not global. 
[... which is not surprising, as this problem is NP-hard.]

Getting started: 
- Forgy method: choose $k$ random sample points to be initial $\mu$’s; go to (2). 
- Random partition: randomly assign each sample point to a cluster; go to (1). 
- $k$-means++: like Forgy, but biased distribution. [Each center is chosen with a preference for points far from previous centers.]

[k-means++ is a little more work, but it works well in practice and theory. Forgy seems to be better than random partition, but Wikipedia mentions some variants of $k$-means for which random partition is better.]
For best results, run $k$-means multiple times with random starts.

Clusters found by running 3-means 6 times on the same sample points, each time starting with a different random partition. The algorithm finds three different local minima.

[Why did we choose that particular objective function to minimize? Partly because it is equivalent to minimizing the following function.]

Equivalent objective fn: the within-cluster variation

Find $y$ that minimizes

$$
\sum_{i=1}^{k} \frac{1}{n_i} \sum_{y_j=i, y_m=i} \|X_j - X_m\|^2
$$

[At the minimizer, this objective function is equal to twice the previous one. It’s a worthwhile exercise to show that—it’s harder than it looks. The nice thing about this expression is that it doesn’t include the means; it’s a function purely of the input points and the clusters we assign them to. So it’s more convincing.]

Normalize the data? [before applying $k$-means]

Same advice as for PCA. Sometimes yes, sometimes no.

[If some features are much larger than others, they will tend to dominate the Euclidean distance. So if you have features in different units of measurement, you probably should normalize them. If you have features in the same unit of measurement, you usually shouldn’t, but it depends on context.]
**k-Medoids Clustering**

Generalizes $k$-means beyond Euclidean distance. [Means aren’t optimal for other distance metrics.]

Specify a distance fn $d(x, y)$ between points $x, y$, aka dissimilarity.

Can be arbitrary; ideally satisfies triangle inequality $d(x, y) \leq d(x, z) + d(z, y)$.

[Sometimes people use the $\ell_1$ norm or the $\ell_\infty$ norm. Sometimes people specify a matrix of pairwise distances between the input points.]

[Suppose you have a database that tells you how many of each product each customer bought. You’d like to cluster together customers who buy similar products for market analysis. But if you cluster customers by Euclidean distance, you’ll get a big cluster of all the customers who have only ever bought one thing. So Euclidean distance is not a good measure of dissimilarity. Instead, it makes more sense to treat each customer as a vector and measure the angle between two customers. If there’s a large angle between customers, they’re dissimilar.]

Replace mean with medoid, the sample point that minimizes total distance to other points in same cluster.

[So the medoid of a cluster is always one of the input points.]

[One difficulty with $k$-means is that you have to choose the number $k$ of clusters before you start, and there isn’t any reliable way to guess how many clusters will best fit the data. The next method, hierarchical clustering, has the advantage in that respect. By the way, there is a whole Wikipedia article on “Determining the number of clusters in a data set.”]

**Hierarchical Clustering**

Creates a tree; every subtree is a cluster. [So some clusters contain smaller clusters.]

Bottom-up, aka agglomerative clustering:

start with each point a cluster; repeatedly fuse pairs.

Top-down, aka divisive clustering:

start with all pts in one cluster; repeatedly split it.

[When the input is a point set, agglomerative clustering is used much more in practice than divisive clustering. But when the input is a graph, it’s the other way around: divisive clustering is more common.]

We need a distance fn for clusters $A, B$:

- **complete linkage:** $d(A, B) = \max\{d(w, x) : w \in A, x \in B\}$
- **single linkage:** $d(A, B) = \min\{d(w, x) : w \in A, x \in B\}$
- **average linkage:** $d(A, B) = \frac{1}{|A||B|} \sum_{w \in A} \sum_{x \in B} d(w, x)$
- **centroid linkage:** $d(A, B) = d(\mu_A, \mu_B)$ where $\mu_S$ is mean of $S$

[The first three of these linkages work for any distance function, even if the input is just a matrix of distances between all pairs of points. The centroid linkage only really makes sense if we’re using the Euclidean distance. But there’s a variation of the centroid linkage that uses the medoids instead of the means, and medoids are defined for any distance function. Moreover, medoids are more robust to outliers than means.]

Greedy agglomerative alg.:

Repeatedly fuse the two clusters that minimize $d(A, B)$

Naively takes $O(n^2)$ time.
[But for complete and single linkage, there are more sophisticated algorithms called CLINK and SLINK, which run in $O(n^2)$ time. A package called ELKI has publicly available implementations.]

Dendrogram: Illustration of the cluster hierarchy (tree) in which the vertical axis encodes all the linkage distances.

Cut dendrogram into clusters by horizontal line according to your choice of # of clusters OR intercluster distance.

[It’s important to be aware that the horizontal axis of a dendrogram has no meaning. You could swap some treenode’s left subtree and right subtree and it would still be the same dendrogram. It doesn’t mean anything that two leaves happen to be next to each other.]
Comparison of average, complete (max), and single (min) linkages. Observe that the complete linkage gives the best-balanced dendrogram, whereas the single linkage gives a very unbalanced dendrogram that is sensitive to outliers (especially near the top of the dendrogram).

[Probably the worst of these is the single linkage, because it's very sensitive to outliers. Notice that if you cut this example into three clusters, two of them have only one sample point. It also tends to give you a very unbalanced tree.]

[The complete linkage tends to be the best balanced, because when a cluster gets large, the furthest point in the cluster is always far away. So large clusters are more resistant to growth than small ones. If balanced clusters are your goal, this is your best choice.]

[In most applications you probably want the average or complete linkage.]

Warning: centroid linkage can cause inversions where a parent cluster is fused at a lower height than its children.

[So statisticians don’t like it, but nevertheless, centroid linkage is popular in genomics.]

[As a final note, all the clustering algorithms we’ve studied so far are unstable, in the sense that deleting a few input points can sometimes give you very different results. But these unstable heuristics are still the most commonly used clustering algorithms. And it’s not clear to me whether a truly stable clustering algorithm is even possible.]
THE GEOMETRY OF HIGH-DIMENSIONAL SPACES

[High-dimensional geometry sometimes acts in ways that are completely counterintuitive, defying our intuitions from low-dimensional geometry.]

Consider a random point $p \sim N(0, I) \in \mathbb{R}^d$. What is the distribution of its length?

[Looking at the one-dimensional normal distribution, you would expect it to be very common that the length is close to zero, a bit less common that the length is close to 1 or −1, and not rare for the length to be close to 2 or −2. But in high dimensions, that intuition is completely wrong.]

If the dimension is very high, the vast majority of the random points are at approximately the same distance from the mean. So they lie in a thin shell. Why? To answer that, let’s study the square of the distance. By Pythagoras’ Theorem, the squared distance from $p$ to the mean is

$$\|p\|^2 = p_1^2 + p_2^2 + \ldots + p_d^2$$

Each component $p_i$ is sampled independently from a univariate normal distribution with mean zero and variance one. The square of a component, $p_i^2$, is said to come from a chi-squared distribution. So is $\|p\|^2$.

$$p_i \sim N(0, 1), \quad p_i^2 \sim \chi^2(1), \quad E[p_i^2] = 1, \quad \text{Var}(p_i^2) = 2, \quad \|p\|^2 \sim \chi^2(d)$$

[Recall that when you add $d$ independent, identically distributed random numbers, you scale their mean and variance by $d$, and the standard deviation is the square root of the variance.]

$$E[\|p\|^2] = d \cdot E[p_i^2] = d$$

$$\text{Var}(\|p\|^2) = d \cdot \text{Var}(p_i^2) = 2d$$

$$\text{SD}(\|p\|^2) = \sqrt{2d}$$

For large $d$, $\|p\|$ is concentrated in a thin shell around radius $\sqrt{d}$ with a thickness proportional to $\sqrt{2d}$. [The mean value of $\|p\|$ isn’t exactly $\sqrt{d}$, but it is close, because the mean of $\|p\|^2$ is $d$ and the standard deviation is much, much smaller. Likewise, the standard deviation of $\|p\|$ isn’t exactly $\sqrt{2d}$, but it’s close.]

[So if $d$ is about a million, imagine a million-dimensional egg whose radius is about 1,000, and the thickness of the shell is about 67, which is about 10 times the standard deviation. The vast majority of random points are in the eggshell. Not inside the egg; actually in the shell itself. It is counterintuitive that random vectors sampled from a high-dimensional normal distribution almost all have almost the same length.]

[There is a statistical principle hiding here. Suppose you want to estimate the mean of a distribution—in this case, the distribution $\chi^2(1)$. The standard way to do that is to sample very many numbers from the distribution and take their mean. The more numbers you sample, the more accurate your estimate is—that is, the smaller the standard deviation of your sample mean is. When we sample a vector from a million-dimensional normal distribution and compute its length, that’s exactly what we’re doing!]

normal.pdf [A one-dimensional normal distribution.]
What about a uniform distribution? Consider concentric spheres of radii \( r \) & \( r - \epsilon \).

![Concentric balls. In high dimensions, almost every point chosen uniformly at random in the outer ball lies outside the inner ball.]

Volume of outer ball \( \propto r^d \)

Volume of inner ball \( \propto (r - \epsilon)^d \)

Ratio of inner ball volume to outer =

\[
\frac{(r - \epsilon)^d}{r^d} = \left(1 - \frac{\epsilon}{r}\right)^d \approx \exp\left(-\frac{\epsilon d}{r}\right) \quad \text{which is small for large } d.
\]

E.g., if \( \frac{\epsilon}{r} = 0.1 \) & \( d = 100 \), inner ball has \( 0.9^{100} = 0.0027\% \) of volume.

Random points from uniform distribution in ball: nearly all are in thin outer shell.

"Gaussian": nearly all are in some thin shell.

Consequences:

- In high dimensions, sometimes the nearest neighbor and 1,000th-nearest neighbor don’t differ much!
- \( k \)-means clustering and nearest neighbor classifiers are less effective for large \( d \).

**Angles between Random Vectors**

What is the angle \( \theta \) between a random \( p \sim \mathcal{N}(0, I) \in \mathbb{R}^d \) and an arbitrary \( q \in \mathbb{R}^d \)?

Without loss of generality, set \( q = [1 \ 0 \ 0 \ldots 0]^T \).

[The value of \( q \) doesn’t matter, because the direction that \( p \) points in is uniformly distributed over all possible directions. By a formula we learned early this semester, the angle between \( p \) and \( q \) is \( \theta \), where . . .]

\[
\cos \theta = \frac{p \cdot q}{\|p\| \|q\|} = \frac{p_1}{\|p\|}
\]

\[
E[\cos \theta] = 0; \quad \text{SD}(\cos \theta) \approx \frac{\text{SD}(p_1)}{\sqrt{d}} = \frac{1}{\sqrt{d}}
\]

If \( d \) is large, \( \cos \theta \) is almost always very close to zero; \( \theta \) is almost always very close to 90°!

[In high-dimensional spaces, two random vectors are almost always very close to orthogonal. To put it another way, an arbitrary vector is almost orthogonal to the vast majority of all the other vectors!]

[A former CS 189/289A head TA, Marc Khoury, has a nice short essay entitled “Counterintuitive Properties of High Dimensional Space”, which you can read at https://marckhoury.github.io/blog/counterintuitive-properties-of-high-dimensional-space]
RANDOM PROJECTION

An alternative to PCA as preprocess for clustering, classification, regression. Approximately preserves distances between points!

[We project onto a random subspace instead of the PCA subspace, but sometimes it preserves distances better than PCA. Because it roughly preserves the distances, algorithms like k-means clustering and nearest neighbor classifiers will give similar results to what they would give in high dimensions, but they run much faster. It works best when you project a very high-dimensional space to a medium-dimensional space.]

Pick a small \( \epsilon \), a small \( \delta \), and a random subspace \( S \subset \mathbb{R}^d \) of dimension \( k \), where

\[
k = \left\lfloor \frac{2 \ln(1/\delta)}{\epsilon^2/2 - \epsilon^3/3} \right\rfloor.
\]

For any pt \( q \), let \( \hat{q} \) be orthogonal projection of \( q \) onto \( S \), multiplied by \( \sqrt{d/k} \).

[The multiplication by \( \sqrt{d/k} \) helps preserve the distances between points after you project.]

Johnson–Lindenstrauss Lemma (modified):
For any two pts \( q, w \in \mathbb{R}^d \), \( (1 - \epsilon) \|q - w\|^2 \leq \|\hat{q} - \hat{w}\|^2 \leq (1 + \epsilon) \|q - w\|^2 \) with probability \( \geq 1 - 2\delta \).

Typical values: \( \epsilon \in [0.02, 0.5] \), \( \delta \in [1/n^3, 0.05] \). [You choose \( \epsilon \) and \( \delta \) according to your needs.]

[The squared distance between two points after projecting changes by less than 2%, or less than 50%, as you wish. In practice, experiment with \( k \) to find the best speed-accuracy tradeoff. If you want all inter-sample-point distances to be accurate, you should set \( \delta \) smaller than \( 1/n^2 \), so you need a subspace of dimension \( \Theta(\log n) \). Reducing \( \delta \) doesn’t cost much (because of the logarithm), but reducing \( \epsilon \) costs more. You can bring 1,000,000 sample points down to a 10,000-dimensional space with at most a 6% error in the distances.] [What is remarkable about this result is that the dimension \( d \) of the input points doesn’t matter!]

[Why does this work? A random projection of \( q - w \) is like taking a random vector and selecting \( k \) components. The mean of the squares of those \( k \) components approximates the mean for the whole population.]

[How do you get a uniformly distributed random projection direction? You can choose each component from a univariate Gaussian distribution, then normalize the vector to unit length. How do you get a random subspace? You can choose \( k \) random directions, then use Gram–Schmidt orthogonalization to make them mutually orthonormal. Interestingly, Indyk and Motwani show that if you skip the expensive normalization and Gram–Schmidt steps, random projection still works almost as well, because random vectors in a high-dimensional space are nearly equal in length and nearly orthogonal to each other with high probability.]
THE PSEUDOINVERSE AND THE SVD

[We’re done with unsupervised learning. For the rest of the semester, we go back to supervised learning.]

[The singular value decomposition can give us insight into the pseudoinverse and its use in least-squares linear regression. If you attended Discussion Section 6, you worked through an explanation of this, but now that I’ve introduced the SVD in Lecture 21, I’d like to summarize it.]

[Let’s understand the pseudoinverse of a diagonal matrix first, then the pseudoinverse of a matrix in general.]

Let $D$ be a diagonal $n \times d$ matrix. [Not necessarily square!]

Find its pseudoinverse $D^+$ by transposing $D$ and replacing every nonzero entry with its reciprocal.

E.g., $D = \begin{bmatrix} 2 & 0 & 0 \\ 0 & 0 & 0 \\ 0 & 0 & 1/3 \\ 0 & 0 & 0 \end{bmatrix}$, $D^+ = \begin{bmatrix} 1/2 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 \\ 0 & 0 & 3 & 0 \\ 0 & 0 & 0 & 0 \end{bmatrix}$, $DD^+ = \begin{bmatrix} 1 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 0 \end{bmatrix}$, $D^+D = \begin{bmatrix} 1 & 0 & 0 \\ 0 & 0 & 0 \\ 0 & 0 & 1 \end{bmatrix}$.

[If $D$ were a square matrix with no zeros on the diagonal, then $D^+$ would be the inverse of $D$, and $DD^+$ and $D^+D$ would be the identity matrix. In general, $DD^+$ and $D^+D$ are always diagonal matrices with 0’s and 1’s only.]

Observe that $DD^+D = D$ and $D^+DD^+ = D^+$. If $D$ is square, $D^2D^+ = D$.

[Because the 0’s and 1’s in $DD^+$ line up with the 0’s and nonzeros in $D$ and $D^+$. This is as close to an “inverse” as a rank-deficient matrix can get. Now let’s consider the pseudoinverse of an arbitrary matrix.]

Let $X$ be any $n \times d$ matrix. Let $X = UDV^\top$ be its SVD. Recall that rank $D = \text{rank } X$.

The Moore–Penrose pseudoinverse of $X$ is $X^+ = VD^+U^\top$.

Observe:

1. $XX^+ = UDV^\top VD^+U^\top = U(DD^+)U^\top$ is symmetric & positive semidefinite.

2. $X^+X = VD^+U^\top UD V^\top = V(D^+D)V^\top$ is symmetric & PSD too.

3. All have the same rank: $D$, $D^+$, $DD^+$, $D^+D$, $X$, $X^+$, $XX^+$, $X^+X$.

4. If $X$ has rank $n$, then $XX^+ = I_{n \times n}$ and $X^+$ is a right inverse.

5. If $X$ has rank $d$, then $X^+X = I_{d \times d}$ and $X^+$ is a left inverse.

6. $XX^+X = X$. Proof: $XX^+X = U(DD^+)U^\top UD V^\top = U(DD^+)V^\top = UD V^\top = X$.

7. $X^+XX^+ = X^+$. [The proof is symmetric to the previous one.]

[Now let’s visualize what the pseudoinverse and the SVD are doing.]

rowcol.pdf [The singular vectors are perpendicular, but we are viewing the planes from oblique angles.]
[If we think of \( X \) as a linear function that maps row \( X \) to col \( X \), and we ignore the null spaces of \( X \) and \( X^\top \), then that linear function is a bijection whose inverse is the pseudoinverse \( X^+ \).]

Linear function \( f : \text{row} \ X \to \text{col} \ X \), \( p \mapsto Xp \) is a bijection. Its inverse is \( f^{-1} : q \mapsto X^+q \).

Let \( r = \text{rank} \ X \). The \( r \) right singular vectors \( v_i \) with \( \delta_i \neq 0 \) are an orthonormal basis for row \( X \). The \( r \) left singular vectors \( u_i \) with \( \delta_i \neq 0 \) are an orthonormal basis for col \( X \).

\( Xv_i = \delta_i u_i \), \( X^+u_i = \frac{1}{\delta_i} v_i \).

[\( X \) maps each right singular vector to some scalar multiple of the corresponding left singular vector. The corresponding singular value tells us how much longer the vector gets when we map it. \( X^+ \) maps each left singular vector to some scalar multiple of a right singular vector.]

[Normally we don’t think of \( X \) as a function from row space to column space. Usually we think of \( X \) as a function from some bigger space \( \mathbb{R}^d \) to a bigger space \( \mathbb{R}^n \). In that case, \( X \) represents a linear function that maps \( \mathbb{R}^d \) to col \( X \), and its pseudoinverse \( X^+ \) maps \( \mathbb{R}^n \) to row \( X \), but in general they’re not bijections. In our figure above, \( X \) might be a \( 4 \times 3 \) matrix, but its rank is only two. So \( X \) maps every point in \( \mathbb{R}^3 \) to a point on the plane on the left. When you map a three-dimensional space down to two dimensions, it can’t be a bijection, so \( X \) doesn’t have an inverse. Just a pseudoinverse.]

[You can think of \( X \) as a function that orthogonally projects a three-dimensional point down onto the row space of \( X \), then uses the bijection above to finish the mapping. Symmetrically, you can think of \( X^+ \) as a function that orthogonally projects a four-dimensional point down onto the column space of \( X \), then uses the inverse bijection. Here’s an illustration of mapping \( p \) to \( Xp \) and \( q \) to \( X^+q \).]

[Now, we can show that the pseudoinverse always gives a good solution in least-squares linear regression, even when \( X^\top X \) is singular.]

Theorem: A solution to the normal equations \( X^\top Xw = X^\top y \) is \( w = X^+y \).

Proof: \( X^\top Xw = X^\top XX^+y = VDU^\top U(DD^+)U^\top y = V(D^2D^+)U^\top y = VDU^\top y = X^\top y \).

If the normal eq’ns have multiple solutions, \( w = X^+y \) is the least-norm solution: i.e., it minimizes \( ||w|| \) among all solutions. [If you attended Discussion Section 6, you might have proven this yourself.]

[This way of solving the normal equations is very helpful when \( X^\top X \) is singular because \( n < d \) or the sample points lie on a subspace of the feature space. But observe that if \( X \) has a very small singular value, the reciprocal of that singular value will be very large and have a very large effect on \( w \); but when that singular value is exactly zero, it has no effect on \( w \)! So when we have a really tiny singular value, should we pretend it is zero? Ridge regression implements this policy to some degree; review Discussion Worksheet 12 for details.]
23 Learning Theory

LEARNING THEORY: WHAT IS GENERALIZATION?

[One thing humans do well is generalize. When you were a young child, you only had to see a few examples of cows before you learned to recognize cows, including cows you had never seen before. You didn’t have to see every cow. You didn’t even have to see log n of the cows.]

[Learning theory tries to explain how machine learning algorithms generalize, so they can classify data they’ve never seen before. It also tries to derive mathematically how much training data we need to generalize well. Learning theory starts with the observation that if we want to generalize, it helps to constrain what hypotheses we allow our learner to consider.]

A range space (aka set system) is a pair \((P, H)\), where

\(P\) is set of all possible test/training points (can be infinite)

\(H\) is hypothesis class, a set of hypotheses (aka ranges, aka classifiers):

- each hypothesis is a subset \(h \subseteq P\) that specifies which points \(h\) predicts are in class C.
  
  [So each hypothesis \(h\) is a 2-class classifier, and \(H\) is a set of sets of points.]

Examples:

1. Power set classifier: \(P\) is a set of \(k\) numbers; \(H\) is the power set of \(P\), containing all \(2^k\) subsets of \(P\).
   
   e.g., \(P = \{1, 2\}, H = \{\emptyset, \{1\}, \{2\}, \{1, 2\}\}\)

2. Linear classifier: \(P = \mathbb{R}^d\); \(H\) is the set of all halfspaces; each halfspace has the form \(\{x : w \cdot x \geq -\alpha\}\).
   
   [In this example, both \(P\) and \(H\) are infinite. In particular, \(H\) contains every possible halfspace—that is, every possible linear classifier in \(d\) dimensions.]

[The power set classifier sounds very powerful, because it can learn every possible hypothesis. But the reality is that it can’t generalize at all. Imagine we have three training points and three test points in a row.]

\[?
  \text{C}
  \text{C}
  ?
  \text{N}
  ?
\]

[The power set classifier can classify these three test points any way you like. Unfortunately, that means it has learned nothing about the test points from the training points. By contrast, the linear classifier can learn only two hypotheses that fit this training data. The leftmost test point must be classified class C, and the rightmost test point must be classified class Not-C. Only the test point in the middle can swing either way. So the linear classifier has a big advantage: it can generalize from a few training points. That’s also a big disadvantage if the data isn’t close to linearly separable, but that’s another story.]

[Now we will investigate how well the training error predicts the test error, and how that differs for these two classifiers.]

Suppose all training pts & test pts are drawn independently from same prob. distribution \(\mathcal{D}\) defined on domain \(P\). [\(\mathcal{D}\) also determines each point’s label. Classes C and Not-C may have overlapping distributions.]

Let \(h \in H\) be a hypothesis [a classifier]. \(h\) predicts a pt \(x\) is in class C if \(x \in h\).

The risk aka generalization error \(R(h)\) of \(h\) is the probability that \(h\) misclassifies a random pt \(x\) drawn from \(\mathcal{D}\)—i.e., the prob. that \(x \in C\) but \(x \notin h\) or vice versa.

[Risk is almost the same as the test error. To be precise, the risk is the mean test error for test points drawn randomly from \(\mathcal{D}\). For a particular test set, sometimes the test error is higher, sometimes lower, but on average it is \(R(h)\). If you had an infinite amount of test data, the risk and the test error would be the same.]
Let $X \subseteq P$ be a set of $n$ training pts drawn from $D$.
The empirical risk aka training error $\hat{R}(h)$ is $\%$ of $X$ misclassified by $h$.
[This matches the definition of empirical risk I gave you in Lecture 12, if you use the 0-1 loss function.]

$h$ misclassifies each training pt w/ prob. $R(h)$, so total misclassified has a binomial distribution.
As $n \to \infty$, $\hat{R}(h)$ better approximates $R(h)$.

[Consider a hypothesis whose risk of misclassification is 25%.
Plotted are probability mass functions of the number of misclassified training points for 20 points and 500 points, respectively. For 20 points, the training error is not a reliable estimate of the risk: the hypothesis might get “lucky” with misleadingly low training error.]

[Hoeffding’s inequality tells us prob. of bad estimate:

$$\Pr(|\hat{R}(h) - R(h)| > \epsilon) \leq 2e^{-2\epsilon^2 n}.$$

[Hoeffding’s inequality is a standard result about how likely it is that a number drawn from a binomial distribution will be far from its mean. If $n$ is big enough, then it’s very unlikely.]
Idea for learning alg: choose \( \hat{h} \in H \) that minimizes \( \hat{R}(\hat{h}) \)! Empirical risk minimization.

[None of the classification algorithms we’ve studied actually do this, but only because it’s computationally infeasible to pick the best hypothesis. Support vector machines can find a linear classifier with zero training error when the training data is linearly separable. But when it isn’t, SVMs try to find a linear classifier with low training error, but they don’t generally find the one with minimum training error. That’s NP-hard.]

[Nevertheless, for the sake of understanding learning theory, we’re going to pretend that we have the computational power to try every hypothesis and pick the one with the lowest training error.]

Problem: if too many hypotheses, some \( h \) with high \( R(h) \) will get lucky and have very low \( \hat{R}(\hat{h}) \)!

[This brings us to a central idea of learning theory. You might think that the ideal learning algorithm would have the largest class of hypotheses, so it could find the perfect one to fit the data. But the reality is that you can have so many hypotheses that some of them just get lucky and score far lower training error than their actual risk. That’s another way to understand what “overfitting” is.]

[More precisely, the problem isn’t too many hypotheses. Usually we have infinitely many hypotheses, and that’s okay. The problem is too many dichotomies.]

**Dichotomies**

A dichotomy of \( X \) is \( X \cap h \), where \( h \in H \).

[A dichotomy picks out the training points that \( h \) predicts are in class C. Think of each dichotomy as a function assigning each training point to class C or class Not-C.]

\[
\begin{array}{ccc}
C & C & N \\
C & N & N \\
C & C & C \\
\end{array}
\]

[Draw this by hand. [dichotomies.pdf]]

[Three examples of dichotomies for three points in a hypothesis class of linear classifiers, and one example (right) that is not a dichotomy.]

[For \( n \) training points, there could be up to \( 2^n \) dichotomies. The more dichotomies there are, the more likely it is that one of them will get lucky and have misleadingly low empirical risk.]

Extreme case: if \( H \) allows all \( 2^n \) possible dichotomies, \( \hat{R}(\hat{h}) = 0 \) even if every \( h \in H \) has high risk.

[If our hypothesis class permits all \( 2^n \) possible assignments of the \( n \) training points to classes, then one of them will have zero training error. But that’s true even if all of the hypotheses are terrible and have a large risk. Because the hypothesis class imposes no structure, we overfit the training points.]

If \( H \) induces \( \Pi \) dichotomies, \( \Pr(\text{at least one dichotomy has } |\hat{R} - R| > \epsilon) \leq \delta, \) where \( \delta = 2\Pi e^{-2\epsilon^2 n}. \)

[Let’s fix a value of \( \delta \) and solve for \( \epsilon \).] Hence with prob. \( \geq 1 - \delta \), for every \( h \in H \),

\[
|\hat{R}(h) - R(h)| \leq \epsilon = \sqrt{\frac{1}{2n} \ln \frac{2\Pi}{\delta}}.
\]

[This tells us that the smaller we make \( \Pi \), the number of possible dichotomies, and the larger we make \( n \), the number of training points, the more accurately the training error will approximate how well the classifier performs on test data.]

smaller \( \Pi \) or larger \( n \) \( \Rightarrow \) training error probably closer to true risk (\& test error).
Smaller $\Pi$ means we’re less likely to overfit. We have less variance, but more bias. This doesn’t necessarily mean the risk will be small. If our hypothesis class $H$ doesn’t fit the data well, both the training error and the test error will be large. In an ideal world, we want a hypothesis class that fits the data well, yet doesn’t produce many dichotomies.

Let $h^* \in H$ minimize $R(h^*)$; “best” classifier.

[Remember we picked the classifier $\hat{h}$ that minimizes the empirical risk. We really want the classifier $h^*$ that minimizes the actual risk, but we can’t know what $h^*$ is. But if $\Pi$ is small and $n$ is large, the hypothesis $\hat{h}$ we have chosen is probably nearly as good as $h^*$.]

With prob. $\geq 1 - \delta$, our chosen $\hat{h}$ has nearly optimal risk:

$$R(\hat{h}) \leq \hat{R}(\hat{h}) + \epsilon \leq \hat{R}(h^*) + \epsilon \leq R(h^*) + 2\epsilon, \quad \epsilon = \sqrt{\frac{1}{2n} \ln \frac{2\Pi}{\delta}}.$$

[This is excellent news! It means that with enough training data and a limit on the number of dichotomies, empirical risk minimization usually chooses a classifier close to the best one in the hypothesis class.]

Choose a $\delta$ and an $\epsilon$.

The sample complexity is the # of training pts needed to achieve this $\epsilon$ with prob. $\geq 1 - \delta$:

$$n \geq \frac{1}{2\epsilon^2} \ln \frac{2\Pi}{\delta}.$$

[If $\Pi$ is small, we won’t need too many training points to choose a good classifier. Unfortunately, if $\Pi = 2^n$ we lose, because this inequality says that $n$ has to be bigger than $n$. So the power set classifier can’t learn much or generalize at all. We need to severely reduce $\Pi$, the number of possible dichotomies. One way to do that is to use a linear classifier.]

The Shatter Function & Linear Classifiers

[How many ways can you divide $n$ points into two classes with a hyperplane?]

\[
\begin{align*}
\text{# of dichotomies: } \Pi_H(X) &= |\{X \cap h : h \in H\}| \quad \epsilon \in [1, 2^n] \text{ where } n = |X| \\
\text{shatter function: } \Pi_H(n) &= \max_{|X|=n, X \subseteq P} \Pi_H(X) \quad \text{[The most dichotomies of any point set of size $n$]}
\end{align*}
\]

Example: Linear classifiers in plane. $H = \text{set of all halfplanes}$. $\Pi_H(3) = 8$:

```
N
C C
C N C
C C C
```

[Draw this by hand. [shatter.pdf] [Linear classifiers can induce all eight dichotomies of these three points. The other four dichotomies are the complements of these four.]
$\Pi_H(4) = 14$:  
[Instead of showing you all 14 dichotomies, let me show you dichotomies that halfplanes cannot learn, which illustrate why no four points have 16 dichotomies.]

[Draw this by hand. ⟨unshatter.pdf⟩] [Examples of dichotomies of four points in the plane that no linear classifier can induce.]

[This isn’t a proof that 14 is the maximum, because we have to show that 15 is not possible for any four points in the plane. The standard proof uses a famous result called Radon’s Theorem.]

Fact: for all range spaces, either $\Pi_H(n)$ is polynomial in $n$, or $\Pi_H(n) = 2^n \forall n \geq 0$.

[This is a surprising fact with deep implications. Imagine that you have $n$ points, some of them training points and some of them test points. Either a range space permits every possible dichotomy of the points, and the training points don’t help you classify the test points at all; or the range space permits only a polynomial subset of the $2^n$ possible dichotomies, so once you have labeled the training points, you have usually cut down the number of ways you can classify the test points dramatically. No shatter function ever occupies the no-man’s-land between polynomial and $2^n$.]

[For linear classifiers, we know exactly how many dichotomies there can be.]

Cover’s Theorem [1965]: linear classifiers in $\mathbb{R}^d$ allow up to $\Pi_H(n) = 2 \sum_{i=0}^{d} \binom{n-1}{i}$ dichotomies of $n$ pts.

For $n \leq d + 1$, $\Pi_H(n) = 2^n$.

For $n \geq d + 1$, $\Pi_H(n) \leq 2 \left(\frac{e(n-1)}{d}\right)^d$  
[Observe that this is polynomial in $n$! With exponent $d$.]

and the sample complexity needed to achieve $R(\hat{h}) \leq \hat{R}(\hat{h}) + \epsilon \leq R(h^*) + 2\epsilon$ with prob. $\geq 1 - \delta$ is

$$n \geq \frac{1}{2\epsilon^2} \left(d \ln \frac{n-1}{d} + d + \ln \frac{4}{\delta}\right).$$  
[Observe that the logarithm turned the exponent $d$ into a factor!]

Corollary: linear classifiers need only $n \in \Theta(d)$ training pts for training error to accurately predict risk or test error.

[In a $d$-dimensional feature space, we need more than $d$ training points to train an accurate linear classifier. But it’s reassuring to know that the number we need is linear in $d$. By contrast, if we have a classifier that permits all $2^n$ possible dichotomies however large $n$ is, then no amount of training data will guarantee that the training error of the hypothesis we choose approximates the true risk.]

[The constant hidden in that big-$\Theta$ notation can be quite large. For example, if you choose $\epsilon = 0.1$ and $\delta = 0.1$, then setting $n = 550d$ will always suffice. (For very large $d$, $n = 342d$ will do.) If you want a lot of confidence that you’ve chosen one of the best hypotheses, you have to pay for it with a large sample size.]

[This sample complexity applies even if you add polynomial features or other features, but you have to count the extra features in $d$. So the number of training points you need increases with the number of polynomial terms.]
VC Dimension

The Vapnik–Chervonenkis dimension of \((P, H)\) is

\[
VC(H) = \max\{n : \Pi_H(n) = 2^n\}. \quad \Leftarrow \text{Can be } \infty.
\]

Say that \(H\) shatters a set \(X\) of \(n\) pts if \(\Pi_H(X) = 2^n\).

VC\((H)\) is size of largest \(X\) that \(H\) can shatter.

[This means that \(X\) is a point set for which all \(2^n\) dichotomies are possible.]

[I told you earlier that if the shatter function isn’t \(2^n\) for all \(n\), then it’s a polynomial in \(n\). The VC dimension is motivated by an observation that sometimes makes it easy to bound that polynomial.]

Theorem: \(\Pi_H(n) \leq \sum_{i=0}^{VC(H)} \binom{n}{i}\). Hence for \(n \geq VC(H)\), \(\Pi_H(n) \leq \left(\frac{en}{VC(H)}\right)^{VC(H)}\).

[So the VC dimension is an upper bound on the exponent of the polynomial. This theorem is useful because often we can find an easy upper bound on the VC dimension. You just need to show that for some number \(n\), no set of \(n\) points can have all \(2^n\) dichotomies.]

Corollary: \(O(VC(H))\) training pts suffice for accuracy. [Again, the hidden constant is big.]

[If the VC dimension is finite, it tells us how the sample complexity grows with the number of features. If the VC dimension is infinite, no amount of training data will make the classifier generalize well.]

Example: Linear classifiers in plane.
Recall \(\Pi_H(3) = 8\): there exist 3 pts shattered by halfplanes.
But \(\Pi_H(4) = 14\): no 4 pts are shattered.

Hence:

- \(VC(H) = 3\) [The VC dimension of halfplanes is 3.]
- \(\Pi_H(n) \leq \frac{e^3}{27} n^3\) [The shatter function is polynomial.]
- \(O(1)\) sample complexity.

[The VC dimension doesn’t always give us the tightest bound. In this example, the VC dimension promises that the number of ways halfplanes can classify the points is at worst cubic in \(n\); but Cover’s Theorem says it’s quadratic in \(n\). In general, linear classifiers in \(d\) dimensions have VC dimension \(d + 1\), which is one dimension looser than the exponent Thomas Cover proved. That’s not a big deal, though, as the sample complexity and the accuracy bound are both based on the logarithm of the shatter function. So if we get the exponent wrong, it only changes a constant in the sample complexity.]

[The important thing is simply to show that there is some polynomial bound on the shatter function at all. VC dimension is not the only way to do that, but often it’s the easiest.]

[The main point you should take from this lecture is that if you want to have generalization, you need to limit the expressiveness of your hypothesis class so that you limit the number of possible dichotomies of a point set. This may or may not increase the bias, but if you don’t limit the number of dichotomies at all, the overfitting could be very bad. If you limit the hypothesis class, your artificial child will only need to look at \(O(d)\) cows to learn the concept of cows. If you don’t, your artificial child will need to look at every cow in the world, and every non-cow too.]
24  Boosting; Nearest Neighbor Classification

ADABOOST (Yoav Freund and Robert Schapire, 1997)

AdaBoost (“adaptive boosting”) is an ensemble method for classification (or regression) that
– reduces bias [compare with random forests, which use ensembles to reduce variance];
– trains learners on weighted sample points [like bagging];
– uses different weights for each learner;
– increases weights of misclassified training points;
– gives bigger votes to more accurate learners.

Input: $n \times d$ design matrix $X$, vector of labels $y \in \mathbb{R}^n$ with $y_i = \pm 1$.

Ideas:
– Train $T$ classifiers $G_1, \ldots, G_T$. [“$T$” stands for “trees”]
– Weight for training point $X_i$ in $G_t$ grows according to how many of $G_1, \ldots, G_{t-1}$ misclassified it.
  [Moreover, if $X_i$ is misclassified by very accurate learners, its weight grows even more.]
  [And, the weight shrinks every time $X_i$ is correctly classified.]
– Train $G_t$ to try harder to correctly classify training pts with larger weights.
– Metalearner is a linear combination of learners. For test point $z$, $M(z) = \sum_{t=1}^{T} \beta_t G_t(z)$.
  Each $G_t$ is $\pm 1$, but $M$ is continuous. Return sign of $M(z)$.

[Remember that in the previous lecture on ensemble methods, I talked briefly about how to assign different weights to training points. It varies for different learning algorithms. For example, in regression we usually modify the risk function by multiplying each point’s loss function by its weight. In a soft-margin support vector machine, we modify the objective function by multiplying each point’s slack by its weight.

[Boosting works with most learning algorithms, but it was originally developed for decision trees, and boosted decision trees are very popular and successful. To weight points in decision trees, we use a weighted entropy where instead of computing the proportion of points in each class, we compute the proportion of weight in each class.]]

In iteration $T$, what classifier $G_T$ and coefficient $\beta_T$ should we choose? Pick a loss fn $L$(prediction, label).

Find $G_T$ & $\beta_T$ that minimize

$$\text{Risk} = \frac{1}{n} \sum_{i=1}^{n} L(M(X_i), y_i), \quad M(X_i) = \sum_{t=1}^{T} \beta_t G_t(X_i).$$

AdaBoost metalearner uses exponential loss function

$$L(\rho, \ell) = e^{-\rho \ell} = \begin{cases} 
  e^{-\rho} & \ell = +1 \\
  e^{\rho} & \ell = -1 
\end{cases}$$

[This loss function is for the metalearner only. We will discover later that the ideal cost function for each individual learner $G_t$ is just the total weight of the misclassified points. In practice, our individual learners are often classification algorithms like decision trees that don’t explicitly try to minimize any loss function at all. Even when the practice doesn’t match the theory, boosting still usually works quite well.]

Important: label $\ell$ is binary, $G_t$ is binary, but $\rho = M(X_i)$ is continuous!

[The exponential loss function has the advantage that it pushes hard against badly misclassified points. So it’s usually better than the squared error loss function for classification in a metalearner. It’s similar to why in neural networks we often prefer the cross-entropy loss function to the squared error.]
Now, we have derived the optimal metalearner!

G must be classified wrongly. Moreover, it’s NP-hard to find the optimal linear classifier, so in practice that’s not always possible. If we’re using a linear classifier on data that’s not linearly separable, some points will be an approximate best learner, not the true minimizer of training error. But that’s okay.

To choose $\beta$, I’ll explain later.

Boosting is usually used with short, imperfect decision trees instead of tall, pure decision trees, for reasons recursive formulation is a nice benefit of choosing the exponential loss function. Notice that a weight shrinks if the point was classified correctly by learner $T$, and grows if the point was misclassified. Moreover, it’s NP-hard to find the optimal linear classifier, so in practice $G_T$ will be an approximate best learner, not the true minimizer of training error. But that’s okay.

You might ask, if we use decision trees, can’t we get zero training error? Usually we can. But interestingly, boosting is usually used with short, imperfect decision trees instead of tall, pure decision trees, for reasons I’ll explain later.

Now, let’s derive the optimal value of $\beta_T$.

To choose $\beta_T$, set $\frac{d}{d\beta_T} \text{Risk} = 0$. 

\[
0 = -\frac{d}{d\beta_T} e^{-\beta_T} \sum_{i=1}^{n} w_i^{(T)} + (e^{\beta_T} + e^{-\beta_T}) \sum_{y_i \neq G_T(X_i)} w_i^{(T)}; \quad \text{[now divide both sides by the first term]}
\]

\[
0 = -1 + (e^{2\beta_T} + 1) \text{err}_T, \quad \text{where err}_T = \frac{\sum_{y_i \neq G_T(X_i)} w_i^{(T)}}{\sum_{i=1}^{n} w_i^{(T)}}; \quad \Leftrightarrow \text{$G_T$’s weighted error rate}
\]

\[
\beta_T = \frac{1}{2} \ln \left( \frac{1 - \text{err}_T}{\text{err}_T} \right).
\]

[So now we have derived the optimal metalearner!]
– If $\text{err}_T = 0$, $\beta_T = \infty$. [So a perfect learner gets an infinite vote.]
– If $\text{err}_T = 1/2$, $\beta_T = 0$. [So a learner with 50% weighted training error gets no vote at all.]

[More accurate learners get bigger votes in the metalearner. Interestingly, a learner with training error worse than 50% gets a negative vote. A learner with 60% error is just as useful as a learner with 40% error; the metalearner just reverses the signs of its votes.]

[Now we can state the AdaBoost algorithm.]

AdaBoost alg:
1. Initialize weights $w_i \leftarrow \frac{1}{n}, \forall i \in [1,n]$.
2. for $t \leftarrow 1$ to $T$
   a. Train $G_t$ with weights $w_i$.
   b. Compute weighted error rate $\text{err} \leftarrow \frac{\sum_{i} \text{misclassified } w_i}{\sum_{i} w_i}$; coefficient $\beta_t \leftarrow \frac{1}{2} \ln \left( \frac{1 - \text{err}}{\text{err}} \right)$.
   c. Reweight pts: $w_i \leftarrow w_i \cdot \begin{cases} e^{\beta_t}, & G_t \text{ misclassifies } X_i \\ e^{-\beta_t}, & \text{otherwise} \end{cases} = w_i \cdot \begin{cases} \frac{1-\text{err}}{\text{err}}, & \text{err} > \frac{1}{2} \\ \frac{\text{err}}{1-\text{err}}, & \text{err} < \frac{1}{2} \end{cases}$
3. return metalearner $h(z) = \text{sign} \left( \sum_{t=1}^{T} \beta_t G_t(z) \right)$.

Why boost decision trees? [As opposed to other learning algorithms?]

– **Boosting reduces bias reliably, but not always variance.** AdaBoost trees are impure to reduce overfitting. [Recall again that random forests use ensembles to reduce variance, but AdaBoost uses them to reduce bias. The AdaBoost variance is more complicated: it often decreases at first, because successive trees focus on different features, but often it later increases. Sometimes boosting overfits after many iterations, and sometimes it doesn’t; it’s hard to predict when it will and when it won’t.]
– Fast. [We’re training many learners, and running many learners at classification time too. Short decision trees that only look at a few features are very fast at both training and testing.]
– No hyperparameter search needed. [Unlike SVMs, neural nets, etc.] [UC Berkeley’s Leo Breiman called AdaBoost with decision trees “the best off-the-shelf classifier in the world.”]
– Easy to make a tree beat 45% training error [or some other threshold] consistently.
– AdaBoost + short trees is a form of subset selection. 
[Features that don’t improve the metalearner’s predictive power enough aren’t used at all. This helps reduce overfitting and running time, especially if there are a lot of irrelevant features.]
– Linear decision boundaries don’t boost well. 
[It takes a lot of boosting to make linear classifiers model really nonlinear decision boundaries well, so SVMs aren’t a great choice. Recall from Discussion Section 9 that ensembles of depth-one stumps are an even worse choice, because they can’t even do XOR. Methods with nonlinear decision boundaries benefit more from boosting, because they allow boosting to reduce the bias faster much. Even depth-two decision trees boost substantially better than depth-one decision trees.]

More about AdaBoost:
– Posterior prob. can be approximated: \( P(Y = 1|x) \approx 1/(1 + e^{-2M(x)}) \).
– Exponential loss is vulnerable to outliers; for corrupted data, use other loss. 
[Loss functions have been derived for dealing with outliers. Unfortunately, they have more complicated weight computations.]
– If every learner beats error \( \mu \) for \( \mu < 50\% \), metalearner training error will eventually be zero. [You will prove this in Homework 7.]
– [The AdaBoost paper and its authors, Freund and Schapire, won the 2003 Gödel Prize, a prize for outstanding papers in theoretical computer science.]

trainboost.pdf, testboost.pdf (ESL, Figures 10.2, 10.3) [Training and testing errors for AdaBoost with stumps, depth-one decision trees that make only one decision each. At left, observe that the training error eventually drops to zero, and even after that the average loss (which is continuous, not binary) continues to decay exponentially. At right, the test error drops to 5.8% after 400 iterations, even though each learner has an error rate of about 46%. AdaBoost with more than 25 stumps outperforms a single 244-node decision tree. In this example no overfitting is observed, but there are other datasets for which overfitting is a problem.]
NEAREST NEIGHBOR CLASSIFICATION

[I saved the simplest classifier for the end of the semester.]

Idea: Given query point \( q \), find the \( k \) training pts nearest \( q \).

Distance metric of your choice.

Regression: Return average label of the \( k \) pts.

Classification: Return class with the most votes from the \( k \) pts OR return histogram of class probabilities.

[The histogram of class probabilities tries to estimate the posterior probabilities of the classes. Obviously, the histogram has limited precision. If \( k = 3 \), then the only probabilities you’ll ever return are 0, 1/3, 2/3, or 1. You can improve the precision by making \( k \) larger, but you might underfit. The histogram works best when you have a huge amount of data.]

<table>
<thead>
<tr>
<th>KNN: K=1</th>
<th>KNN: K=10</th>
<th>KNN: K=100</th>
</tr>
</thead>
</table>

[Examples of 1-NN, 10-NN, and 100-NN. A larger \( k \) smooths out the boundary. In this example, the 1-NN classifier is badly overfitting the data, and the 100-NN classifier is badly underfitting. The 10-NN classifier does well: it’s reasonably close to the Bayes decision boundary (purple). Generally, the ideal \( k \) depends on how dense your data is. As your data gets denser, the best \( k \) increases.]

[There are theorems showing that if you have a lot of data, nearest neighbors can work quite well.]

Theorem (Cover & Hart, 1967):

As \( n \to \infty \), the 1-NN error rate is \( < 2B - B^2 \) where \( B = \text{Bayes risk} \).

if only 2 classes, \( \leq 2B - 2B^2 \)

[There are a few technical requirements of this theorem. The most important is that the training points and the test points all have to be drawn independently from the same probability distribution—just like in our last lecture, on learning theory. Here, we are using the 0-1 loss to define the Bayes risk; so the Bayes risk is the smallest possible error rate over that distribution. The theorem applies to any separable metric space, so it’s not just for the Euclidean metric.]

[By the way, this Cover is the same Thomas Cover who proved Cover’s Theorem from the last lecture. He was a professor in Electrical Engineering and Statistics at Stanford until his death in 2012, and these are the first and third journal articles he published.]

Theorem (Fix & Hodges, 1951):

As \( n \to \infty \), \( k \to \infty \), \( k/n \to 0 \), \( k\)-NN error rate converges to \( B \). [Which means Bayes optimal.]
25 Nearest Neighbor Algorithms: Voronoi Diagrams and k-d Trees

NEAREST NEIGHBOR ALGORITHMS

Exhaustive k-NN Alg.

Given query point \( q \):
- Scan through all \( n \) sample pts, computing (squared) distances to \( q \).
- Maintain a max-heap with the \( k \) shortest distances seen so far.
  [Whenever you encounter a sample point closer to \( q \) than the point at the top of the heap, you remove the heap-top point and insert the better point. Obviously you don’t need a heap if \( k = 1 \) or even 5, but if \( k = 101 \) a heap will substantially speed up keeping track of the \( k \)th-shortest distance.]

Time to train classifier: 0 [This is the only \( O(0) \)-time algorithm we’ll learn this semester.]
Query time: \( O(nd + n \log k) \)
  expected \( O(nd + k \log n \log k) \) if random pt order

[It’s a cute theoretical observation that you can slightly improve the expected running time by randomizing the point order so that only expected \( O(k \log n) \) heap operations occur. But in practice I can’t recommend it; you’ll probably lose more from cache misses than you’ll gain from fewer heap operations.]

Can we preprocess training pts to obtain sublinear query time?

2–5 dimensions: Voronoi diagrams
Medium dim (up to \( \sim 30 \)): k-d trees
Large dim: exhaustive k-NN, but can use PCA or random projection
  locality sensitive hashing [still researchy, not widely adopted]

Voronoi Diagrams

Let \( X \) be a point set. The Voronoi cell of \( w \in X \) is
\[
\text{Vor} = \{ p \in \mathbb{R}^d : \| p - w \| \leq \| p - v \| \ \forall v \in X \}
\]
[A Voronoi cell is always a convex polyhedron or polytope.]
The Voronoi diagram of \( X \) is the set of \( X \)'s Voronoi cells.
Voronoi diagrams sometimes arise in nature (salt flats, giraffe, crystallography).

Believe it or not, the first published Voronoi diagram dates back to 1644, in the book “Principia Philosophiae” by the famous mathematician and philosopher René Descartes. He claimed that the solar system consists of vortices. In each region, matter is revolving around one of the fixed stars (vortex.pdf). His physics was wrong, but his idea of dividing space into polyhedral regions has survived.

Size (e.g., # of vertices) $\in O(n^{d/2})$

[This upper bound is tight when $d$ is a small constant. As $d$ grows, the tightest asymptotic upper bound is somewhat smaller than this, but the complexity still grows exponentially with $d$.]

...but often in practice it is $O(n)$.

[Here I’m leaving out a constant that may grow exponentially with $d$.]
Point location: Given query point \( q \in \mathbb{R}^d \), find the point \( w \in X \) for which \( q \in \text{Vor} \ w \).

[We need a second data structure that can perform this search on a Voronoi diagram efficiently.]

2D: \( O(n \log n) \) time to compute \( \text{V.d.} \) and a trapezoidal map for pt location

\[ O(\log n) \] query time \quad [because of the trapezoidal map]

[That’s a pretty great running time compared to the linear query time of exhaustive search.]

\( d \)D: Use binary space partition tree (BSP tree) for pt location

[Unfortunately, it’s difficult to characterize the running time of this strategy, although it is likely to be reasonably fast in 3–5 dimensions.]

1-NN only!

[A standard Voronoi diagram supports only 1-nearest neighbor queries. If you want the \( k \) nearest neighbors, there is something called an order-\( k \) Voronoi diagram that has a cell for each possible \( k \) nearest neighbors.

But nobody uses those, for two reasons. First, the size of an order-\( k \) Voronoi diagram is \( O(k^2 n) \) in 2D, and worse in higher dimensions. Second, there’s no software available to compute one.]

[There are also Voronoi diagrams for other distance metrics, like the \( L_1 \) and \( L_\infty \) norms.]

[Voronoi diagrams are good for 1-nearest neighbor queries in 2 or 3 dimensions, maybe 4 or 5, but \( k \)-d trees are much simpler and probably faster in 6 or more dimensions.]

\( k \)-d Trees

“Decision trees” for NN search. Differences: \quad [compared to decision trees]

– Choose splitting feature w/greatest width: feature \( i \) in \( \max_{j,k}(X_{ji} - X_{ki}) \).

[With nearest neighbor search, we don’t care about the entropy. Instead, what we want is that if we draw a sphere around the query point, it won’t intersect very many boxes of the decision tree. So it helps if the boxes are nearly cubical, rather than long and thin.]

Cheap alternative: rotate through the features. \quad [We split on the first feature at depth 1, the second feature at depth 2, and so on. This builds the tree faster, by a factor of \( O(d) \).]

– Choose splitting value: median point for feature \( i \); OR \( \frac{X_{ji} + X_{ki}}{2} \).

Median guarantees \( \lfloor \log_2 n \rfloor \) tree depth; \( O(nd \log n) \) tree-building time.

[... or just \( O(n \log n) \) time if you rotate through the features. An alternative to the median is splitting at the box center, which improves the aspect ratios of the boxes, but it could unbalance your tree. A compromise strategy is to alternate between medians at odd depths and centers at even depths, which also guarantees an \( O(\log n) \) depth.]

– Each internal node stores a sample point. \quad [... that lies in the node’s box. Usually the splitting point.]

[Some \( k \)-d tree implementations have points only at the leaves, but it’s better to have points in internal nodes too, so when we search the tree, we often stop searching before we reach a leaf.]

root represents \( \mathbb{R}^2 \)

right halfplane

lower right
quarter plane

[Draw this by hand. kdtreestructure.pdf]
Goal: given query pt $q$, find a sample pt $w$ such that $\|q - w\| \leq (1 + \epsilon) \|q - s\|$, where $s$ is the closest sample pt.

$\epsilon = 0 \Rightarrow$ exact NN; \hspace{1cm} $\epsilon > 0 \Rightarrow$ approximate NN.

Query alg. maintains:
- Nearest neighbor found so far (or $k$ nearest).
- Binary min-heap of unexplored subtrees, keyed by distance from $q$.

[Draw this by hand.  \[kdtreequery.pdf\]  [A query in progress.]

[Each subtree represents an axis-aligned box. The query tries to avoid searching most of the boxes/subtrees by searching the boxes close to $q$ first. We measure the distance from $q$ to a box and use it as a key for the subtree in the heap. The search stops when the distance from $q$ to the $k$th-nearest neighbor found so far \\(\leq\) the distance from $q$ to the nearest unexplored box (times $1 + \epsilon$). For example, in the figure above, the query never visits the box at far lower right, because it doesn’t intersect the circle.]

Alg. for 1-NN query:

$$Q \leftarrow \text{heap containing root node with key zero}$$

$$r \leftarrow \infty$$

while $Q$ not empty and $(1 + \epsilon) \cdot \text{minkey}(Q) < r$

$$B \leftarrow \text{removemin}(Q)$$

$$w \leftarrow B’s \text{ sample point}$$

$$r \leftarrow \min\{r, \text{dist}(q, w)\}$$

[For speed, store square of $r$ instead.]

$$B’, B” \leftarrow \text{child boxes of } B$$

if $(1 + \epsilon) \cdot \text{dist}(q, B’) < r$ then insert($Q, B’, \text{dist}(q, B’)$) [The key for $B’$ is dist($q, B’$)]

if $(1 + \epsilon) \cdot \text{dist}(q, B”) < r$ then insert($Q, B”, \text{dist}(q, B”)$)

return point that determined $r$

For $k$-NN, replace “$r$” with a max-heap holding the $k$ nearest neighbors
[...just like in the exhaustive search algorithm.]

Works with any $L_p$ norm for $p \in [1, \infty]$. [k-d trees are not limited to the Euclidean ($L_2$) norm.]

Why $\epsilon$-approximate NN?

[Draw this by hand.  \[kdtreeproblem.pdf\]  [A worst-case exact NN query.]
[In the worst case, we may have to visit every node in the k-d tree to find the exact nearest neighbor. In that case, the k-d tree is slower than simple exhaustive search. This is an example where an approximate nearest neighbor search can be much faster. In practice, settling for an approximate nearest neighbor sometimes improves the speed by a factor of 10 or even 100, because you don’t need to look at most of the tree to do a query. This is especially true in high dimensions—remember that in high-dimensional space, the nearest point often isn’t much closer than a lot of other points.] 

Software: ANN (U. Maryland), FLANN (U. British Columbia), GeRaF (U. Athens) [random forests!]

Example: im2gps

[I want to emphasize the fact that exhaustive nearest neighbor search really is one of the first classifiers you should try in practice, even if it seems too simple. So here’s an example of a modern research paper that uses 1-NN and 120-NN search to solve a problem.] 

Paper by James Hays and [our own] Prof. Alexei Efros. 
[Goal: given a query photograph, determine where on the planet the photo was taken. Called geolocalization. They evaluated both 1-NN and 120-NN. What they did not do, however, is treat each photograph as one long vector. That’s okay for tiny digits, but too expensive for millions of travel photographs. Instead, they reduced each photo to a small descriptor made up of a variety of features that extract the essence of each photo.] 

[Bottom line: With 120-NN, their most sophisticated implementation came within 64 km of the correct location about 50% of the time.]

RELATED CLASSES [if you like machine learning, consider taking these courses]

CS 182 (spring, fall?): Deep Neural Networks 
CS C281A (fall?): Statistical Learning Theory [C281A is the most direct continuation of CS 189/289A.] 
EECS 127/227AT (both), EECS 227BT (fall): Numerical Optimization [a core part of ML] 
[It’s hard to overemphasize the importance of numerical optimization to machine learning, as well as other CS fields like graphics, theory, and scientific computing.] 
EECS 126 (both): Random Processes [Markov chains, expectation maximization, PageRank] 
EE C106A/B (fall/spring): Intro to Robotics [dynamics, control, sensing] 
Math 110: Linear Algebra [but the real gold is in Math 221] 
Math 221 (fall): Matrix Computations [how to solve linear systems, compute SVDs, eigenvectors, etc.] 
CS 194-26/294-26 (fall): Computational Photography (Efros, Kanazawa) 
CS C267 (spring): Scientific Computing [parallelization, practical matrix algebra, some graph partitioning] 
CS C280 (spring): Computer Vision 
CS 285 (fall): Deep Reinforcement Learning 
CS 288 (fall): Natural Language Processing 
CS 294-43 (spring): Visual Object/Activity Recognition (Darrell) 
CS 294-150 (spring): ML & Biology (Listgarten) 
CS 294-158 (spring): Deep Unsupervised Learning (Abbeel) 
CS 294-162 (fall?): ML Systems (Gonzalez) 
CS 294-173 (spring): Learning for 3D Vision (Kanazawa) 
CS 294-182 (spring): ML & Econ (Haghtalab) 
VS 265: Neural Computation
A  Bonus Lecture: Spectral Graph Clustering

SPECTRAL GRAPH CLUSTERING

Input:  Weighted, undirected graph \( G = (V, E) \). No self-edges.
\( w_{ij} = \) weight of edge \( (i, j) = (j, i) \); zero if \( (i, j) \notin E \).

[Think of the edge weights as a similarity measure. A big weight means that the two vertices want to be in the same cluster. So the circumstances are the opposite of the last lecture on clustering. Then, we had a distance or dissimilarity function, so small numbers meant that points wanted to stay together. Today, big numbers mean that vertices want to stay together.]

Goal: Cut \( G \) into 2 (or more) pieces \( G_i \) of similar sizes,
but don’t cut too much edge weight.
[That’s a vague goal. There are many ways to make this precise.
Here’s a typical goal, which we’ll solve approximately.]
e.g., Minimize the sparsity \( \frac{\text{Cut}(G_1, G_2)}{\text{Mass}(G_1) \cdot \text{Mass}(G_2)} \), aka cut ratio
where \( \text{Cut}(G_1, G_2) = \) total weight of cut edges
\( \text{Mass}(G_1) = \) # of vertices in \( G_1 \) OR assign masses to vertices

[The denominator “Mass(\( G_1 \)) Mass(\( G_2 \))” penalizes imbalanced cuts.]

\[\text{graph.pdf}\]  [Four cuts. All edges have weight 1.
Upper left: the minimum bisection; a bisection is perfectly balanced.
Upper right: the minimum cut. Usually very unbalanced; not what we want.
Lower left: the sparsest cut, which is good for many applications.
Lower right: the maximum cut; in this case also the maximum bisection.]

Sparsest cut, min bisection, max cut all NP-hard.
[Today we will look for an approximate solution to the sparsest cut problem.]
[We will turn this combinatorial graph cutting problem into algebra.]
Let \( n = |V| \). Let \( y \in \mathbb{R}^n \) be an indicator vector:

\[
y_i = \begin{cases} 
1 & \text{vertex } i \in G_1, \\
-1 & \text{vertex } i \in G_2.
\end{cases}
\]

Then \( w_{ij} \frac{(y_i - y_j)^2}{4} = \begin{cases} 
w_{ij} & \text{if } (i, j) \text{ is cut,} \\
0 & \text{if } (i, j) \text{ is not cut.}
\end{cases} \]

\[
\text{Cut}(G_1, G_2) = \sum_{(i,j) \in E} w_{ij} \frac{(y_i - y_j)^2}{4} \quad \text{[This is quadratic, so let’s try to write it with a matrix.]} \\
= \frac{1}{4} \sum_{(i,j) \in E} \left( w_{ij} y_i^2 - 2w_{ij} y_i y_j + w_{ij} y_j^2 \right) \\
= \frac{1}{4} \left( \sum_{(i,j) \in E} -2w_{ij} y_i y_j + \sum_{i=1}^n y_i^2 \sum_{k \neq i} w_{ik} \right) \\
= \frac{y^T Ly}{4},
\]

where \( L_{ij} = \begin{cases} 
-w_{ij}, & i \neq j, \\
\sum_{k \neq i} w_{ik}, & i = j.
\end{cases} \)

\( L \) is symmetric, \( n \times n \) Laplacian matrix for \( G \).

[L is effectively a matrix representation of \( G \). For the purpose of partitioning a graph, there is no need to distinguish edges of weight zero from edges that are not in the graph.]

[We see that minimizing the weight of the cut is equivalent to minimizing the Laplacian quadratic form \( y^T Ly \). This lets us turn graph partitioning into a problem in matrix algebra.]

[Usually we assume there are no negative weights, in which case \( \text{Cut}(G_1, G_2) \) can never be negative, so it follows that \( L \) is positive semidefinite.]

Define \( \mathbf{1} = [1 \quad 1 \quad \ldots \quad 1]^T \); then \( \mathbf{L} \mathbf{1} = \mathbf{0} \), so \( \mathbf{1}^T \mathbf{L} = 0 \), so \( \mathbf{1} \) is an eigenvector of \( L \) with eigenvalue \( 0 \).

[If \( G \) is a connected graph and all the edge weights are positive, then this is the only zero eigenvalue. But if \( G \) is not connected, \( L \) has one zero eigenvalue for each connected component of \( G \). It’s easy to prove, but time prevents me.]
Bisection: exactly \( n/2 \) vertices in \( G_1 \), \( n/2 \) in \( G_2 \). Write \( \mathbf{1}^T y = 0 \).

[So we have reduced graph bisection to this constrained optimization problem.]

Minimum bisection:

\[
\begin{array}{l}
\text{Find } y \text{ that minimizes } y^T Ly \\
\text{subject to } \forall i, y_i = 1 \text{ or } y_i = -1 \\
\text{and } \mathbf{1}^T y = 0
\end{array}
\]

← binary constraint

← balance constraint

Also NP-hard. We relax the binary constraint. → fractional vertices!

[A very common approach in combinatorial optimization algorithms is to relax some of the constraints so a discrete problem becomes a continuous problem. Intuitively, this means that you can put \( 1/3 \) of vertex 7 in graph \( G_1 \) and the other \( 2/3 \) of vertex 7 in graph \( G_2 \). You can even put \( -1/2 \) of vertex 7 in graph \( G_1 \) and \( 3/2 \) of vertex 7 in graph \( G_2 \). This sounds crazy, but the continuous problem is much easier to solve than the combinatorial problem. After we solve the continuous problem, we will round the vertex values to \( +1/-1 \), and we’ll hope that our solution is still close to optimal.]

[But we can’t just drop the binary constraint. We still need some constraint to rule out the solution \( y = 0 \).]

New constraint: \( y \) must lie on hypersphere of radius \( \sqrt{n} \).

[Draw this by hand. circle.pdf] [Instead of constraining \( y \) to lie at a vertex of the hypercube, we constrain \( y \) to lie on the hypersphere through those vertices.]

Relaxed problem:

\[
\begin{array}{l}
\text{Minimize } y^T Ly \\
\text{subject to } y^T y = n \\
\text{and } \mathbf{1}^T y = 0
\end{array}
\]

\[=\]

\[
\text{Minimize } \frac{y^T Ly}{y^T y} = \text{Rayleigh quotient of } L \& y \\
\text{(subject to same two constraints)}
\]

[cylinder.pdf] [The isosurfaces of \( y^T Ly \) are elliptical cylinders. The gray cross-section is the hyperplane \( \mathbf{1}^T y = 0 \). We seek the point that minimizes \( y^T Ly \), subject to the constraints that it lies on the gray cross-section and that it lies on a sphere centered at the origin.]
[You should remember this Rayleigh quotient from the lecture on PCA. As I said then, when you see a Rayleigh quotient, you should smell eigenvectors nearby. The $y$ that minimizes this Rayleigh quotient is the eigenvector with the smallest eigenvalue. We already know what that eigenvector is: it’s $1$. But that violates our balance constraint. As you should recall from PCA, when you’ve used the most extreme eigenvector and you need an orthogonal one, the next-best optimizer of the Rayleigh quotient is the next eigenvector.]

Let $\lambda_2 =$ second-smallest eigenvalue of $L$.

Eigenvector $v_2$ is the Fiedler vector. $v_2$ solves the relaxed problem.

[It would be wonderful if every component of the Fiedler vector was 1 or $-1$, but that happens more or less never. So we round $v_2$. The simplest way is to round all positive entries to 1 and all negative entries to $-1$.

But in both theory and practice, it’s better to choose the threshold as follows.]

Spectral partitioning alg:
- Compute Fiedler vector $v_2$ of $L$
- Round $v_2$ with a sweep cut:
  = Sort components of $v_2$.
  = Try the $n - 1$ cuts between successive components. Choose min-sparsity cut.
    [If we’re clever about updating the sparsity, we can try all these cuts in time linear in the number of edges in $G$.]

[Left: example of a graph partitioned by the sweep cut. Right: what the un-rounded Fiedler vector looks like.]
One consequence of relaxing the binary constraint is that the balance constraint no longer forces an exact bisection. But that’s okay; we’re cool with a slightly unbalanced cut if it means we cut fewer edges. Even though our discrete problem was the minimum bisection problem, our relaxed, continuous problem will be an approximation of the sparsest cut problem. This is a bit counterintuitive.

![A graph for which an unbalanced cut (left) is sparser than a balanced one (right).]

**Vertex Masses**

Sometimes you want the notion of balance to accord more prominence to some vertices than others. We can assign masses to vertices.

Let $M$ be diagonal matrix with vertex masses on diagonal.

New balance constraint: $\mathbf{1}^T My = 0$.

[This new balance constraint says that $G_1$ and $G_2$ should each have the same total mass. It turns out that this new balance constraint is easier to satisfy if we also revise the sphere constraint a little bit.]

New ellipsoid constraint: $y^T My = \text{Mass}(G) = \sum M_{ii}$.

[Instead of a sphere, now we constrain $y$ to lie on an axis-aligned ellipsoid.]

Now solution is Fiedler vector of generalized eigensystem $Lv = \lambda Mv$.

[Most algorithms for computing eigenvectors and eigenvalues of symmetric matrices can easily be adapted to compute eigenvectors and eigenvalues of symmetric generalized eigensystems too.]

[For the grad students, here’s the most important theorem in spectral graph partitioning.]

Fact: Sweep cut finds a cut w/ sparsity $\leq \sqrt{2\lambda_2 \max \frac{Lw}{Mw}}$: Cheeger’s inequality.

The optimal cut has sparsity $\geq \lambda_2/2$.

[So the spectral partitioning algorithm is an approximation algorithm, albeit not one with a constant factor of approximation. Cheeger’s inequality is a very famous result in spectral graph theory, because it’s one of the most important cases where you can relax a combinatorial optimization problem to a continuous optimization problem, round the solution, and still have a provably decent solution to the original combinatorial problem.]
Vibration Analogy

For intuition about spectral partitioning, think of the eigenvectors as vibrational modes in a physical system of springs and masses. Each vertex models a point mass that is constrained to move freely along a vertical rod. Each edge models a vertical spring with rest length zero and stiffness proportional to its weight, pulling two point masses together. The masses are free to oscillate sinusoidally on their rods. The eigenvectors of the generalized eigensystem $L\mathbf{v} = \lambda M\mathbf{v}$ are the vibrational modes of this physical system, and their eigenvalues are proportional to their frequencies.

These illustrations show the first four eigenvectors for two simple graphs. On the left, we see that the first eigenvector is the eigenvector of all 1’s, which represents a vertical translation of all the masses in unison. That’s not really a vibration, which is why the eigenvalue is zero. The second eigenvector is the Fiedler vector, which represents the vibrational mode with the lowest frequency. Each component indicates the amplitude with which the corresponding point mass oscillates. At any point in time as the masses vibrate, roughly half the mass is moving up while half is moving down. So it makes sense to cut between the positive components and the negative components. The third eigenvector also gives us a nice bisection of the grid graph, entirely different from the Fiedler vector. Some more sophisticated graph clustering algorithms use multiple eigenvectors.

I want to emphasize that spectral partitioning takes a global view of a graph. It looks at the whole gestalt of the graph and finds a good cut. By comparison, the clustering algorithms we saw last lecture were much more local in nature, so they’re easier to fool.
Greedy Divisive Clustering

Partition $G$ into 2 subgraphs; recursively partition them. [The sparsity is a good criterion for graph clustering. Use $G$’s sparsest cut to divide it into two subgraphs, then recursively cut them. You can stop when you have the right number of clusters. Alternatively, you can make a finer tree and then prune it back.]

The Normalized Cut

Set vertex $i$’s mass $M_{ii} = L_{ii}$. [Sum of edge weights adjoining vertex $i$.] [That is how we define a normalized cut, which turns out to be a good choice for many different applications.] Popular for image segmentation. [Image segmentation is the problem of looking at a photograph and separating it into different objects. To do that, we define a graph on the pixels.] For pixels with coordinate $p_i$, brightness $b_i$, use graph weights

$$w_{ij} = \exp\left(-\frac{||p_i - p_j||^2}{\alpha} - \frac{|b_i - b_j|^2}{\beta}\right) \quad \text{or zero if } ||p_i - p_j|| \text{ large.}$$

[We choose a distance threshold, typically less than 4 to 10 pixels apart. Pixels that are far from each other aren’t connected. $\alpha$ and $\beta$ are empirically chosen constants. It often makes sense to choose $\beta$ proportional to the variance of the brightness values.]
Invented by [our own] Prof. Jitendra Malik and his student Jianbo Shi.
**B  Bonus Lecture: Multiple Eigenvectors; Latent Factor Analysis**

**Clustering w/Multiple Eigenvectors**

[When we use the Fiedler vector for spectral graph clustering, it tells us how to divide a graph into two graphs. If we want more than two clusters, we can use divisive clustering: we repeatedly cut the subgraphs into smaller subgraphs by computing their Fiedler vectors. However, there are several other methods to subdivide a graph into \( k \) clusters in one shot that use multiple eigenvectors rather than just the Fiedler vector \( v_2 \). These methods sometimes give better results. They use \( k \) eigenvectors in a natural way to cluster a graph into \( k \) subgraphs.]

For \( k \) clusters, compute first \( k \) eigenvectors \( v_1, v_2, \ldots, v_k \) of generalized eigensystem \( Lv = \lambda Mv \).

Scale them so that \( v_j^\top M v_j = 1 \). E.g., \( v_1 = \frac{1}{\sqrt{\sum_i M_{ii}}} \). Now \( V^\top M V = I \). [The eigenvectors are \( M \)-orthogonal.]

\[
V = \begin{bmatrix}
    v_1 \\
    \vdots \\
    v_k \\
\end{bmatrix}
\]

\[
V = \begin{bmatrix}
    V_1 \\
    \vdots \\
    V_n \\
\end{bmatrix}
\]

[V’s columns are the eigenvectors with the \( k \) smallest eigenvalues.]

[Yes, we do include the all-1’s vector \( v_1 \) as one of the columns of \( V \).]

Normalizing each row \( V_i \) to unit length.

[Now you can think of the spectral vectors as points on a unit sphere centered at the origin.]

[k-means cluster these vectors.]

[Because all the spectral vectors lie on the sphere, \( k \)-means clustering will cluster together vectors that are separated by small angles.]
[Comparison of point sets clustered by $k$-means—just $k$-means by itself, that is—vs. a spectral method. To create a graph for the spectral method, we use an exponentially decaying function to assign weights to pairs of points, like we used for image segmentation but without the brightnesses.]

Invented by [our own] Prof. Michael Jordan, Andrew Ng [when he was still a student at Berkeley], Yair Weiss.

[This wasn’t the first algorithm to use multiple eigenvectors for spectral clustering, but it has become one of the most popular.]
LATENT FACTOR ANALYSIS [aka Latent Semantic Indexing]

[You can think of this as dimensionality reduction for matrices.]

Suppose $X$ is a term-document matrix: [aka bag-of-words model]
row $i$ represents document $i$; column $j$ represents term $j$. [Term = word.]
Term-document matrices are usually sparse, meaning most entries are zero.
$X_{ij} =$ occurrences of term $j$ in doc $i$
    better: $\log (1 + \text{occurrences})$ [So frequent words don’t dominate.]
Better still is to weight the entries so rare words give big entries and common words like “the” give small entries. To do that, you need to know how frequently each word occurs in general. I’ll omit the details, but this is the common practice.]

Recall $\text{SVD } X = UDV^\top = \sum_{i=1}^{d} \delta_i u_i v_i^\top$. Suppose $\delta_i \leq \delta_j$ for $i \geq j$.

Unlike PCA, we usually don’t center $X$.
For large $\delta_i$, $u_i$ and $v_i$ represent a cluster of documents & terms.
- Large components in $u_i$ mark docs using similar/related terms, i.e., a genre.
- E.g., $u_1$ might have large components for the romance novels,
- $v_1$ ” ” ” ” for terms “passion,” “ravish,” “bodice” …

[… and $\delta_1$ would give us an idea how much bigger the romance novel market is than the markets for every other genre of books.]

[$v_1$ and $u_1$ tell us that there is a large subset of books that tend to use the same large subset of words. We can read off the words by looking at the larger components of $v_1$, and we can read off the books by looking at the larger components of $u_1$.]

[The property of being a romance novel is an example of a latent factor. So is the property of being the sort of word used in romance novels. There’s nothing in $X$ that tells you explicitly that romance novels exist, but the similar vocabulary is a hidden connection between them that gives them a large singular value. The vector $u_1$ reveals which books have that genre, and $v_1$ reveals which words are emphasized in that genre.]

Like clustering, but clusters overlap: if $u_1$ picks out romances & $u_2$ picks out histories, they both pick out historical romances.

[So you can think of latent factor analysis as a sort of clustering that permits clusters to overlap. Another way in which it differs from traditional clustering is that the $u$-vectors contain real numbers, and so some points have stronger cluster membership than others. One book might be just a bit romance, another a lot.]
Application in market research:
identifying consumer types (hipster, suburban mom) & items bought together.
[For applications like this, the first few singular vectors are the most useful. Most of the singular vectors are mostly noise, and they have small singular values to tell you so. This motivates approximating a matrix by using only some of its singular vectors.]

Truncated SVD $X' = \sum_{i=1}^{r} \delta_i u_i v_i^T$ is a low-rank approximation of $X$, of rank $r$.  [Assuming $\delta_r > 0$.]

[We choose the singular vectors with the largest singular values, because they carry the most information.]

\[
\begin{array}{c}
X' \\
\downarrow \\
n \times d
\end{array}
= \begin{array}{c}
\begin{bmatrix}
\delta_1 & \cdots & 0 \\
0 & \cdots & \delta_r \\
\end{bmatrix} \\
r \times r
\end{array}
\begin{array}{c}
\begin{bmatrix}
u_1 \\
\vdots \\
u_r
\end{bmatrix} \\
r \times d
\end{array}
\begin{array}{c}
n \times r
\end{array}
\]

[Draw this by hand.]

$X'$ is the rank-$r$ matrix that minimizes the [squared] Frobenius norm

$\|X - X'\|_F^2 = \sum_{i,j} (X_{ij} - X'_{ij})^2$

Applications:

- Fuzzy search. [Suppose you want to find a document about gasoline prices, but the document you want doesn’t have the word “gasoline”; it has the word “petrol.” One cool thing about the reduced-rank matrix $X'$ is that it will probably associate that document with “gasoline,” because the SVD tends to group synonyms together.]

- Denoising. [The idea is to assume that $X$ is a noisy measurement of some unknown matrix that probably has low rank. If that assumption is partly true, then the reduced-rank matrix $X'$ might be better than the input $X$.]

- Matrix compression. [As you can see above, if we use a low-rank approximation with a small rank $r$, we can express the approximate matrix as an SVD that takes up much less space than the original matrix. Often this low-rank approximation supports faster matrix computations.]

- Collaborative filtering: fills in unknown values, e.g., user ratings.
  [Suppose the rows of $X$ represents Netflix users and the columns represent movies. The entry $X_{ij}$ is the review score that user $i$ gave to movie $j$. But most users haven’t reviewed most movies. We want to fill in the missing values. Just as the rank reduction will associate “petrol” with “gasoline,” it will tend to associate users with similar tastes in movies, so the reduced-rank matrix $X'$ can predict ratings for users who didn’t supply any.]
Big Five (BF) model of personality:

- O: openness
- C: conscientiousness
- E: extraversion
- A: agreeableness
- N: neuroticism

[Researchers have found that these five personality factors are approximately orthogonal to each other. They are highly heritable and highly stable during adulthood.]

Can we predict these traits from 3D faces?

[Studies have shown that people looking at photographs of static faces with neutral expressions can identify the traits better than chance, especially for conscientiousness, extraversion, and agreeableness. This experiment asks whether machine learning can do the same with 3D reconstructions of faces. The subjects were 834 Han Chinese volunteers in Shanghai, China. We don’t know whether any of these results might generalize to people who are not Han Chinese.]

[The faces were scanned in high-resolution 3D and a non-rigid face registration system was used to fit a grid of 32,251 vertices to each face in a manner that maps each vertex to an appropriate landmark on the face. (They call this “anatomical homology.”) So the design matrix $X$ was $834 \times 100,053$, representing 834 subjects with 32,251 3D features each.]

[Subject personalities were evaluated with a self-questionnaire, namely our own Berkeley Personality Lab’s Big Five Inventory, translated into Chinese. The authors treated men and women separately.]
Uses **partial least squares** (PLS) to find associations between personality & faces.

[Everything from here to the end is spoken, not written.]

Partial least squares (PLS) is like a supervised version of PCA. It takes in two matrices $X$ and $Y$ with the same number of rows. In our example, $X$ is the face data and $Y$ is the personality data for the 834 subjects. Like PCA, PLS finds a set of vectors in face space that we think of as the most important components. But whereas PCA looks for the directions of maximum variation in $X$, PLS looks for the directions in $X$ that maximize the correlation with the personality traits in matrix $Y$.

The researchers found the top 20 or so PLS components and used cross-validation to decide which components have predictive power for each personality trait. They found that the top two components for extraversion in women were predictive, but no components for the other four traits in women were predictive. Men are easier to analyze: they found two or three components were predictive for each of extraversion, agreeableness, conscientiousness, and neuroticism in men. However, the correlations were statistically significant only for agreeableness and conscientiousness.

![male.pdf](attachment://male.pdf) [The relationship between male faces, agreeableness, and conscientiousness. The large, colored faces are the mean faces. Colors indicate the values in the most predictive PLS component vector.]

More agreeable men correlate with much wider mouths that look a bit smiley even when neutral; stronger, forward jaws; wider noses; and shorter faces, especially shorter in the forehead, compared to less agreeable men. More conscientious men tend to have higher, wider eyebrows; wider, opened eyes; a withdrawn upper lip with more mouth tension; and taller faces with more pronounced brow ridges (the bone protuberance above the eyes). The authors note that men with low A and C scores look both more relaxed and more indifferent.
More extraverted women correlate with rounder faces, especially in profile, with a more protruding nose and lips but a recessed chin, whereas the introverts have more flat, square-shaped faces. To my eyes, the extraverts also have more expressive mouths.

It’s interesting is that physiognomy, the art of judging character from facial shape, used to be considered a pseudoscience, but it’s been making a comeback in recent years with the help of machine learning. One reason it fell into disrepute is because, historically, it was sometimes applied across races in fallacious and insulting ways. But if you want to train classifiers that guess people’s personalities with some accuracy, you probably need a different classifier for each race. This is a classifier trained exclusively for one race, Han Chinese, which is probably part of why it works as well as it does. If you tried to train one classifier to work on many different races, I suspect its performance would be much worse.

Another thing that’s notable is that the authors were able to find statistically significant correlations for some personality traits, but the majority of traits defeated them. So while physiognomy has some predictive power, it’s only weakly predictive. It’s an open question whether machine learning will ever be able to predict personality from visual information substantially better than this or not. Adding a time dimension and incorporating people’s movements and dynamic facial expressions seems like a promising way to improve personality predictions.

Tools like this raise some ethical issues. The one that concerns me the most is that, if tools like this are emerging now, many governments probably already had similar tools ten years ago, and have probably been using them to profile us.

One student asked whether these methods might be used by employers to screen prospective employees. I think that tools like this are inferior to simply giving an interviewee a personality test. Such tests are legal in the USA, so long as their questions are not found to violate an employee’s right to privacy and the results are not used to discriminate against legally protected groups. The most troubling part of using physiognomy to screen employees would not be that personality testing is unlawful. (It isn’t, and quite a few companies do it.) It would be that physiognomy isn’t nearly accurate enough. An employer who uses a poorly designed or unvalidated personality test to make personnel decisions might run a higher risk that a court might rule that the test could have a discriminatory effect, violating Title VII of the Civil Rights Act of 1964. Also, they probably won’t make good decisions. But perhaps in the future, better measurements, better statistical procedures, and better algorithms might overcome these problems.