CS 480/680 Introduction to Machine Learning

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Introduction

Definition: Artificial Intelligence

A scientific field concerned with the development of algorithms that allow computers to learn without being explicitly programmed.

Definition: Machine Learning

A branch of Artificial Intelligence, which focuses on methods that learn from data and make predictions on unseen data.

There are three phases to machine learning: 1) training; 2) prediction; 3) evaluation

Definition: Paradigms of ML Algorithms

Supervised: learning with labeled data (\mathbf{x}, y) .

Unsupervised: discover patterns in unlabeled data x. Semi-supervised: using both labeled and unlabeled data.

Part I Classical Machine Learning

Perceptron

Dataset: Each column is a data point: n in total; each has d features. y is the label vector.

		Training samples					Test samples	
	\mathbf{x}_1	\mathbf{x}_2	\mathbf{x}_3	\mathbf{x}_4		\mathbf{x}_n	\mathbf{x}_1'	\mathbf{x}_2'
	0	1	0	1		1	1	0.9
$\mathbb{R}^d ightarrow$ Feature $ig\{$	0	0	1	1		0	1	1.1
1 J realure	:	÷	÷	÷	٠.	:	:	:
l	1	0	1	0	• • •	1	1	-0.1
Label y	+	+	_	+		_	?	?

 \mathbf{x}_1' and \mathbf{x}_2' are test samples whose labels need to be predicted.

Definition: Bag of Words Representation of Text

If a word appears, the feature is 1.

Email Spam Detection Example:
$$X = [\mathbf{x}_1, \dots, \mathbf{x}_n] \in \mathbb{R}^{d \times n}$$
 and $y = [y_1, \dots, y_n] \in \{\pm 1\}^n$.

Each column of X is an email $\mathbf{x}_i \in \mathbb{R}^d$, each with d binary features. Each entry in y is a label $y_i \in \{\pm 1\}$ indicating spam or not. Given new email \mathbf{x}' , predict spam or not.

For an OR dataset, the data is linearly separable meaning you can

2.1 Linear Separator

Definition: Inner Product

$$\langle a, b \rangle := \sum_{j} a_{j} b_{j}$$

where a_j and b_j are the jth elements of vectors **a** and **b**.

Definition: Linear Function

 $\forall \alpha, \beta \in \mathbb{R}, \forall \mathbf{x}, \mathbf{z} \in \mathbb{R}^d,$

$$f(\alpha \mathbf{x} + \beta \mathbf{z}) = \alpha f(\mathbf{x}) + \beta f(\mathbf{z})$$

equivalently, $\exists \mathbf{w} \in \mathbb{R}^d$ such that $f(\mathbf{x}) = \langle \mathbf{x}, \mathbf{w} \rangle = \sum_j x_j w_j$.

Proof. (\Longrightarrow) Let $\mathbf{w}=[f(\mathbf{e}_1),\ldots,f(\mathbf{e}_d)]$, where \mathbf{e}_i is the *i*th coordinate vector.

$$f(\mathbf{x}) = f(x_1\mathbf{e}_1 + x_2\mathbf{e}_2 + \dots + x_d\mathbf{e}_d)$$

= $x_1f(\mathbf{e}_1) + \dots + x_df(\mathbf{e}_d) = \langle \mathbf{x}, \mathbf{w} \rangle$

 (\Leftarrow) We have

$$f(\alpha \mathbf{x} + \beta \mathbf{z}) = \langle \alpha \mathbf{x} + \beta \mathbf{z}, \mathbf{w} \rangle$$
$$= \alpha \langle \mathbf{x}, \mathbf{w} \rangle + \beta \langle \mathbf{z}, \mathbf{w} \rangle$$
$$= \alpha f(\mathbf{x}) + \beta f(\mathbf{z})$$

Definition: Affine Function

 $\exists \mathbf{w} \in \mathbb{R}^d, b \in \mathbb{R} \text{ such that } f(\mathbf{x}) = \langle \mathbf{x}, \mathbf{w} \rangle + b.$

Definition: Thresholding

$$\operatorname{sign}(t) = \begin{cases} +1 & t > 0 \\ -1 & t \le 0 \end{cases}$$

Definition: Linear Separator

 \mathbf{w} and b uniquely determine the linear separator.

$$\hat{y} = \operatorname{sign}(\langle \mathbf{x}, \mathbf{w} \rangle + b) = \begin{cases} +1 & \langle \mathbf{x}, \mathbf{w} \rangle + b > 0 \\ -1 & \langle \mathbf{x}, \mathbf{w} \rangle + b \le 0 \end{cases}$$

 \mathbf{w} is orthogonal to the decision boundary $H = {\mathbf{x} : \langle \mathbf{w}, \mathbf{x} \rangle + b = 0}.$

Proof. Any vector with both head and tail in H can be written as $\mathbf{x}\mathbf{x}' = \mathbf{x}' - \mathbf{x}$ for $\mathbf{x}, \mathbf{x}' \in H$.

$$\langle \mathbf{w}, \mathbf{x} \mathbf{x}' \rangle = \langle \mathbf{w}, \mathbf{x}' - \mathbf{x} \rangle$$

$$= \langle \mathbf{w}, \mathbf{x}' \rangle - \langle \mathbf{w}, \mathbf{x} \rangle$$

$$= -b - (-b)$$

$$= 0$$

b does not affect the orthogonality. Holds for any H. The length of \mathbf{w} does not matter in determining the decision boundary.

2.2 Perceptron Algorithm

Algorithm 1 Training Perceptron

Input: Dataset = $(\mathbf{x}_i, y_i) \in \mathbb{R}^d \times \{\pm 1\} : i = 1, ..., n$, initialization $\mathbf{w}_0 \in \mathbb{R}^d$ and $b_0 \in \mathbb{R}$ Output: \mathbf{w} and b (a linear classifier sign($\langle \mathbf{x}, \mathbf{w} \rangle + b$))

for t = 1, 2, ... do

receive index $I_t \in \{1, ..., n\}$ (can be random)

if $y_{I_t}(\langle \mathbf{x}_{I_t}, \mathbf{w} \rangle + b) \leq 0$ then $\mathbf{w} \leftarrow \mathbf{w} + y_{I_t}\mathbf{x}_{I_t}$ $b \leftarrow b + y_{I_t}$

Typically set $\mathbf{w}_0 = \mathbf{0}$ and $b_0 = 0$. $\operatorname{score}_{\mathbf{w},b}(\mathbf{x}) = \hat{y} = \langle \mathbf{w}, \mathbf{x} \rangle + b$.

Perceptron as an Optimization Problem

Find
$$\mathbf{w} \in \mathbb{R}^d$$
, $b \in \mathbb{R}$ such that $\forall i, y_i(\langle \mathbf{x}_i, \mathbf{w} \rangle + b) > 0$

When a mistake happens on (\mathbf{x}, y) we want to update $\mathbf{w}_{k+1} \leftarrow \mathbf{w}_k + y\mathbf{x}$ and $b_{k+1} = b_k + y$:

$$y[\langle \mathbf{x}, \mathbf{w}_{k+1} \rangle + b_{k+1}] = y[\langle \mathbf{x}, \mathbf{w}_k + y\mathbf{x} \rangle + b_k + y] = y[\langle \mathbf{x}, \mathbf{w}_k \rangle + b_k] + \underbrace{\|\mathbf{x}\|_2^2 + 1}_{\text{always positive}}$$

A trick to hide the bias term of the affine function is to pad a constant 1 to the end of each \mathbf{x} :

$$\langle \mathbf{x}, \mathbf{w} \rangle + b = \left\langle \underbrace{\begin{pmatrix} \mathbf{x} \\ 1 \end{pmatrix}}_{\mathbf{x}_{pad}}, \underbrace{\begin{pmatrix} \mathbf{w} \\ b \end{pmatrix}}_{\mathbf{w}_{pad}} \right\rangle$$

The update when a mistake happens on $\langle \mathbf{x}, y \rangle$ is

$$\begin{cases} \mathbf{w} \leftarrow \mathbf{w} + y\mathbf{x} \\ b \leftarrow b + y \end{cases} \iff \mathbf{w}_{pad} \leftarrow \mathbf{w}_{pad} + y\mathbf{x}_{pad}$$

2.3 Convergence

Theorem Convergence (Linearly Separable Case); Block (1962); Novikoff (1962)

Suppose $\exists \mathbf{w}^*$ such that $y_i \langle \mathbf{x}_i, \mathbf{w}^* \rangle > 0$ for all i. Assume that $\|\mathbf{x}_i\|_2 \leq C$ for all i and we normalize the \mathbf{w}^* such that $\|\mathbf{w}^*\|_2 = 1$. Let us define the margin $\gamma := \min_i |\langle \mathbf{x}_i, \mathbf{w}^* \rangle|$. Then the Perceptron algorithm converges after C^2/γ^2 mistakes.

Remember normalizing \mathbf{w}^* does not change the decision boundary of the linear classifier.

Proof. Recall that the update is $\mathbf{w} \leftarrow \mathbf{w} + y\mathbf{x}$ when a mistake happens on (\mathbf{x}, y) . Consider the effect of an update on $\langle \mathbf{w}, \mathbf{w}^* \rangle$:

$$\langle \mathbf{w} + y\mathbf{x}, \mathbf{w}^* \rangle = \langle \mathbf{w}, \mathbf{w}^* \rangle + y \langle \mathbf{x}, \mathbf{w}^* \rangle = \langle \mathbf{w}, \mathbf{w}^* \rangle + |\langle \mathbf{x}, \mathbf{w}^* \rangle| \ge \langle \mathbf{w}, \mathbf{w}^* \rangle + \gamma$$

This means that for each update, $\langle \mathbf{w}, \mathbf{w}^* \rangle$ grows by at least $\gamma > 0$. Consider the effect of an update on $\langle \mathbf{w}, \mathbf{w} \rangle$:

$$\langle \mathbf{w} + y\mathbf{x}, \mathbf{w} + y\mathbf{x} \rangle = \langle \mathbf{w}, \mathbf{w} \rangle + \underbrace{2y \langle \mathbf{w}, \mathbf{x} \rangle}_{\leq 0} + \underbrace{y^2 \langle \mathbf{x}, \mathbf{x} \rangle}_{\in [0, C^2]} \leq \langle \mathbf{w}, \mathbf{w} \rangle + C^2$$

This means for each update, $\langle \mathbf{w}, \mathbf{w} \rangle$ grows by at most C^2 .

Let $\mathbf{w}_0 = \mathbf{0}$. Now we know that after M updates that $\langle \mathbf{w}, \mathbf{w}^* \rangle \geq M \gamma$ and $\langle \mathbf{w}, \mathbf{w} \rangle \leq M C^2$. Thus,

$$1 \ge \cos(\mathbf{w}, \mathbf{w}^*) = \frac{\langle \mathbf{w}, \mathbf{w}^* \rangle}{\|\mathbf{w}\| \|\mathbf{w}^*\|}$$
$$\ge \frac{M_{\gamma}}{\sqrt{MC^2} \times 1}$$
$$= \sqrt{M} \frac{\gamma}{C}$$

This implies $M \leq C^2/\gamma^2$. The larger the margin γ is, the more linearly separable the data is, and hence the fast the Perceptron algorithm will converge.

We want to minimize the Perceptron loss. Define

$$l(\mathbf{w}, \mathbf{x}_t, y) = -y_t \langle \mathbf{w}, \mathbf{x}_t \rangle \mathbb{I}[\text{mistake on } \mathbf{x}_t] = -\min\{y_t \langle \mathbf{w}, \mathbf{x}_t \rangle, 0\}$$
$$L(\mathbf{w}) = -\frac{1}{n} \sum_{t=1}^{n} y_t \langle \mathbf{w}, \mathbf{x}_t \rangle \mathbb{I}[\text{mistake on } \mathbf{x}_t]$$

where \mathbb{I} is the indicator function.

Definition: (Stochastic) Gradient Descent Update

$$\mathbf{w}_{t+1} = \mathbf{w}_t - \eta_t \nabla_{\mathbf{w}} l(\mathbf{w}_t, \mathbf{x}_t, y_t) = \mathbf{w}_t + \eta_t y_t \mathbf{x}_t \mathbb{I}[\text{mistake on } \mathbf{x}_t]$$

Set step size $\eta_t = 1$. If a mistake on (\mathbf{x}_t, y_t) , then

$$\mathbf{w}_{t+1} = \mathbf{w}_t + y_t \mathbf{x}_t$$

Proposition

Perceptron is not unique since the algorithm stops as long as there is no mistake.

We want to maximize the smallest margin: support vector machines.

There is no separating hyperplane for the XOR dataset.

Proof. Suppose $\exists \mathbf{w}, b \text{ such that } y(\langle \mathbf{w}, \mathbf{x} \rangle + b) > 0.$

•
$$\mathbf{x}_1 = (0,0), y_1 = - \implies b < 0$$

•
$$\mathbf{x}_2 = (1,0), y_2 = + \implies w_1 + b > 0$$

•
$$\mathbf{x}_3 = (0,1), y_3 = + \implies w_2 + b > 0 \implies w_1 + w_2 + 2b > 0$$

•
$$\mathbf{x}_4 = (1,1), y_4 = - \implies w_1 + w_2 + b < 0 \implies b > 0$$

This is a contradiction.

Theorem Minsky and Papert (1969); Block and Levin (1970)

If there is no perfect separating hyperplane, then the Perceptron algorithm cycles.

Definition: Soft-Margin

Induced by a reasonable loss l and regularizer reg

$$\min_{\mathbf{w}} \hat{\mathbb{E}}l(y\hat{y}) + reg(\mathbf{w}), \text{ s.t. } \hat{y} := \langle \mathbf{x}, \mathbf{w} \rangle + b$$

Penalizing a mistake by the loss l, but not infinitely large.

Stopping Perceptron:

- Maximum number of iterations is reached
- Maximum allowed runtime is reached
- Training error stops changing
- Validation error stops decreasing

Multiclass Perceptron: Let c be the total number of classes. There are two methods:

1. One vs. All

Let class k be positive and all other classes negative. Train Perceptron \mathbf{w}_k . There are a total of c imbalanced Perceptrons. Predict according to the highest score:

$$\hat{y} := \arg\max_{k} \langle \mathbf{x}, \mathbf{w}_k \rangle$$

2. One vs. One Let class k be positive and class l be negative, discard other classes. Train Perceptron $\mathbf{w}_{k,l}$; $\binom{c}{2}$ balanced Perceptrons. Predict by majority vote:

$$\hat{y} := \arg\max_{k} \sum_{i \neq k} \langle \mathbf{x}, \mathbf{w}_{k,l} \rangle > 0$$

Linear Regression

3.1 General Regression

Definition: Regression

Given training data (\mathbf{x}_i, y_i) where $\mathbf{x}_i \in \mathcal{X} \subseteq \mathbb{R}^d$ is the feature vector for the *i*-th training sample and $y_i \in \mathcal{Y} \subseteq \mathbb{R}^t$ are the *t* responses where t = 1 or $t = \infty$, find $f : \mathcal{X} \to \mathcal{Y}$ such that

$$f(\mathbf{x}_i) \approx y_i$$

Theorem (Exact Interpolation Always Possible)

For any finite training data (\mathbf{x}_i, y_i) : i = 1, ..., n, there exist infinitely many functions f such that for all i,

$$f(\mathbf{x}_i) = y_i$$

- Cannot decide on a unique f.
- On new data \mathbf{x} , our prediction $\hat{y} = f(\mathbf{x})$ can vary significantly.
- Leveraging prior knowledge of f is important.
- Simplest explanation is usually correct.

Prior knowledge on the functional form of f such as linear or nonlinear.

There are three results of regression: underfitting, good fit/robust, overfitting.

Training and test data are both i.i.d. (independent and identically distributed) samples from the same unknown distribution \mathcal{P} . $((\mathbf{X}_i, Y_i) \sim \mathcal{P} \text{ and } (\mathbf{X}, Y) \sim \mathcal{P})$

Definition: Least Squares Regression

$$\min_{f:\mathcal{X}\to\mathcal{Y}} \mathbb{E} \left\| f(\mathbf{X}) - Y \right\|_2^2$$

where we use l_2 loss to measure error.

The square is to make the calculation of the gradient easier.

Definition: Regression Function

Optimal regression function to the least squares regression:

$$f^*(\mathbf{x}) = m(\mathbf{x}) = \mathbb{E}[Y|\mathbf{X} = \mathbf{x}]$$

Bias-Variance Decomposition: To show the regression function is optimal. The expectation is with respect to both X and Y.

$$\mathbb{E} \|f(\mathbf{X}) - Y\|_{2}^{2} = \mathbb{E} \|f(\mathbf{X}) - m(\mathbf{X}) + m(\mathbf{X}) - Y\|_{2}^{2}$$

$$= \mathbb{E} \|f(\mathbf{X}) - m(\mathbf{X})\|_{2}^{2} + \mathbb{E} \|m(\mathbf{X}) - Y\|_{2}^{2}$$

$$+ \underbrace{2\mathbb{E} \langle f(\mathbf{X}) - m(\mathbf{X}), m(\mathbf{X}) - Y \rangle}_{=0}$$

$$= \mathbb{E} \|f(\mathbf{X}) - m(\mathbf{X})\|_{2}^{2} + \underbrace{\mathbb{E} \|m(\mathbf{X}) - Y\|_{2}^{2}}_{\text{noise (variance)}}$$

Showing the 0 term:

$$\mathbb{E} \langle f(\mathbf{X}) - m(\mathbf{X}), m(\mathbf{X}) - Y \rangle = \mathbb{E}_{\mathbf{X}} [\mathbb{E}_{Y|\mathbf{X}} \langle f(\mathbf{X}) - m(\mathbf{X}), m(\mathbf{X}) - Y \rangle]$$
$$= \mathbb{E}_{\mathbf{X}} \langle f(\mathbf{X}) - m(\mathbf{X}), m(\mathbf{X}) - m(\mathbf{X}) \rangle = 0$$

So the following holds true for any f. The noise variance is a constant term w.r.t. to f. We aim to choose $f \approx m$ to minimize the squared error.

$$\mathbb{E} \|f(\mathbf{X}) - Y\|_2^2 = \mathbb{E} \|f(\mathbf{X}) - m(\mathbf{X})\|_2^2 + \underbrace{\mathbb{E} \|m(\mathbf{X}) - Y\|_2^2}_{\text{poise (variance)}}$$

Remember $m(\mathbf{x}) = \mathbb{E}[Y|\mathbf{X} = \mathbf{x}]$ is our gold rule, however it is inaccessible since we do not know the conditional distribution, we are learning f from the training data D.

Let f_D be the regressor learned on D. We have proved that

$$\mathbb{E}_{\mathbf{X},Y} \|f_D(\mathbf{X}) - Y\|_2^2 = \mathbb{E}_{\mathbf{X}} \|f_D(\mathbf{X}) - m(\mathbf{X})\|_2^2 + \underbrace{\mathbb{E}_{\mathbf{X},Y} \|m(\mathbf{X}) - Y\|_2^2}_{\text{noise (variance)}}$$

Define $\overline{f}(\mathbf{X}) = \mathbb{E}_D[f_D(\mathbf{X})]$. We can break down the first term

$$\mathbb{E}_{D}\mathbb{E}_{\mathbf{X}} \|f_{D}(\mathbf{X}) - m(\mathbf{X})\|_{2}^{2} = \mathbb{E}_{D,\mathbf{X}} \|f_{D}(\mathbf{X}) - \overline{f}(\mathbf{X}) + \overline{f}(\mathbf{X}) - m(\mathbf{X})\|_{2}^{2}$$

$$= \mathbb{E}_{\mathbf{X}} \|\overline{f}(\mathbf{X}) - m(\mathbf{X})\|_{2}^{2} + \mathbb{E}_{D,\mathbf{X}} \|f_{D}(\mathbf{X}) - \overline{f}(\mathbf{X})\|_{2}^{2}$$

$$\underbrace{2\mathbb{E}_{D,\mathbf{X}} \left\langle \overline{f}(\mathbf{X}) - m(\mathbf{X}), f_{D}(\mathbf{X}) - \overline{f}(\mathbf{X}) \right\rangle}_{=0}$$

Note that

$$\mathbb{E}_{D,\mathbf{X}} \left\langle f(\mathbf{X}) - m(\mathbf{X}), f_D(\mathbf{X}) - \overline{f}(\mathbf{X}) \right\rangle = \mathbb{E}_{\mathbf{X}} \mathbb{E}_D \left\langle \overline{f}(\mathbf{X}) - m(\mathbf{X}), f_D(\mathbf{X}) - \overline{f}(\mathbf{X}) \right\rangle$$

$$= \mathbb{E}_{\mathbf{X}} \left\langle \overline{f}(\mathbf{X}) - m(\mathbf{X}), \mathbb{E}_D[f_D(\mathbf{X})] - \overline{f}(\mathbf{X}) \right\rangle$$

$$= \mathbb{E}_{\mathbf{X}} \left\langle \overline{f}(\mathbf{X}) - m(\mathbf{X}), \overline{f}(\mathbf{X}) - \overline{f}(\mathbf{X}) \right\rangle = 0$$

Definition: Bias-Variance Trade-Off

$$\underbrace{\mathbb{E}_{D,\mathbf{X},Y} \|f_D(\mathbf{X} - Y)\|_2^2}_{\text{test error}} = \underbrace{\mathbb{E}_{\mathbf{X}} \|\mathbb{E}_D[f_D(\mathbf{X})] - m(\mathbf{X})\|_2^2}_{\text{bias}^2} + \underbrace{\mathbb{E}_{D,\mathbf{X}} \|f_D(\mathbf{X}) - \mathbb{E}_D[f_D(\mathbf{X})]\|_2^2}_{\text{variance}} + \underbrace{\mathbb{E}_{\mathbf{X},Y} \|m(\mathbf{X}) - Y\|_2^2}_{\text{noise (variance)}}$$

Bias² \downarrow , as the model capacity \uparrow (model is more expressively powerful). Variance \uparrow , as the model capacity \uparrow (prediction of model is less stable).

Replace expectation with sample average $(\mathbf{X}_i, Y_i) \sim P$:

$$\min_{f:\mathcal{X} \to \mathcal{Y}} \hat{\mathbb{E}} \| f(\mathbf{X}) - Y \|_{2}^{2} := \frac{1}{n} \sum_{i=1}^{n} \| f(\mathbf{X}_{i}) - Y_{i} \|_{2}^{2}$$

Uniform law of large numbers: as training data size $n \to \infty$,

$$\hat{\mathbb{E}} \to \mathbb{E}$$
 and (hopefully) $\arg \min \hat{\mathbb{E}} \to \arg \min \mathbb{E}$

3.2 Linear Regression

Recall the affine function $f(\mathbf{x}) = W\mathbf{x} + \mathbf{b}$ with $W \in \mathbb{R}^{t \times d}$ and $\mathbf{b} \in \mathbb{R}^t$ and padding $\mathbf{x} \leftarrow \begin{pmatrix} \mathbf{x} \\ 1 \end{pmatrix}, W \leftarrow [W, \mathbf{b}]$, hence $f(\mathbf{x}) = W\mathbf{x}$.

In matrix form, $X = [\mathbf{x}_1, \dots, \mathbf{x}_n] \in \mathbb{R}^{(d+1)\times n}, Y = [y_1, \dots, y_n] \in \mathbb{R}^{t\times n}$ and define $||A||_F = \sqrt{\sum_{ij} a_{ij}^2}$ where a_{ij} is the the ijth element of A.

$$\frac{1}{n}\sum_{i=1}^{n} \|f(\mathbf{x}_i) - y_i\|_2^2 = \frac{1}{n} \|WX - Y\|_F^2$$

Theorem (Fermat's Necessary Condition for Extremity)

If **w** is a minimizer (or maximizer) of a differentiable function f over an open set, then $f'(\mathbf{w}) = \mathbf{0}$.

Training: The loss function is defined as

$$Loss(W) = \frac{1}{n} \|WX - Y\|_F^2$$

Taking the derivative $\nabla_W \operatorname{Loss}(W) = \frac{2}{n}(WX - Y)X^T$ and setting to zero results in the **normal equation**:

$$WXX^T = YX^T \implies W = YX^T(XX^T)^{-1}$$

Once you solve for W on the training set (X,Y), we can predict on unseen data X_{test} :

$$\hat{Y}_{test} = WX_{test}$$

Evaluating test error if labels were available:

$$\frac{1}{n_{test}} \left\| Y_{test} - \hat{Y}_{test} \right\|_F^2$$

We can compare the predicted value and errors to reduce training error as a means to reduce test error.

Ill-Conditioning: Let $X = \begin{pmatrix} 0 & \epsilon \\ 1 & 1 \end{pmatrix}$ and $y = \begin{pmatrix} 1 & -1 \end{pmatrix}$. If we solve the linear least squares regression:

$$\mathbf{w} = yX^T(XX^T)^{-1} = \begin{pmatrix} 1 & -1 \end{pmatrix} \begin{pmatrix} -1/\epsilon & 1 \\ 1/\epsilon & 0 \end{pmatrix} = \begin{pmatrix} -2/\epsilon & 1 \end{pmatrix}$$

Slight perturbation leads to chaotic behaviour. This happens whenever X is ill-conditioned, i.e. close to rank deficient.

Rank deficient $X \implies$ two columns in X are linearly dependent, but the corresponding y's might be different. This is a contradiction and lead to an unstable \mathbf{w} .

Definition: Ridge Regression

$$\min_{W} \frac{1}{n} \|WX - Y\|_{F}^{2} + \lambda \|W\|_{F}^{2}$$

The normal equation is $W(XX^T + n\lambda I) = YX^T$. The regularization constant λ controls the trade-off. $\lambda = 0$ reduces to ordinary linear regression while $\lambda = \infty$ reduces to $W = \mathbf{0}$.

$$Loss(W) = \frac{1}{n} \|WX - Y\|_F^2 + \lambda \|W\|_F^2$$

$$\nabla_W Loss(W) = \frac{2}{n} (WX - Y)X^T + 2\lambda W = 0$$

$$WXX^T - YX^T + W(\lambda nI) = 0$$

$$W(XX^T + \lambda nI) = YX^T$$

$$W = YX^T (XX^T + \lambda nI)^{-1}$$

Data Augmentation:

$$\frac{1}{n} \|WX - Y\|_F^2 + \lambda \|W\|_F^2 = \frac{1}{n} \left\| W \underbrace{\left(X \quad \sqrt{n\lambda}I \right)}_{\tilde{X}} - \underbrace{\left(Y \quad \mathbf{0} \right)}_{\tilde{Y}} \right\|_F^2$$

Augment X with $\sqrt{n\lambda}I$, i.e. p data points $\mathbf{x}_j = \sqrt{n\lambda}\mathbf{e}_j$ for $j = 1, \dots, p$, augment Y with zero, which shrinks W towards the origin. Therefore, regularization = data augmentation.

Theorem (Lagrangian)

A hard constraint problem

$$\min_{\theta} L(\theta) = \frac{1}{n} \sum_{i=1}^{n} l(\theta; \mathbf{x}_i, y_i), \text{ s.t.} R(\theta) \le r$$

is equivalent to a soft constraint problem

$$\min_{\theta} L_R(\theta) = \frac{1}{n} \sum_{i=1}^{n} l(\theta; \mathbf{x}_i, y_i) + \lambda^* R(\theta)$$

for some regularization parameter $\lambda^* > 0$, in the sense that they have the same solutions.

Sparsity: Regularization \iff constraint:

$$\min_{\|W\|_F \le \gamma} \frac{1}{n} \|WX - Y\|_F^2$$

Ridge regression \rightarrow dense W which is more computation and harder to interpret.

Lasso (Tibshirani 1996): $\min_{\|W\|_1 \le \gamma} \frac{1}{n} \|WX - Y\|_F^2$.

$$\min_{W} \frac{1}{n} \|WX - Y\|_{F}^{2} + \lambda \|W\|_{1}$$

Logistic Regression

Though called logistic regression, it is for linear classification.

Definition: Logit

$$\langle \mathbf{x}, \mathbf{w} \rangle$$

Let $\mathcal{Y} = \{0, 1\}$ and directly learn confidence $p(\mathbf{x}; \mathbf{w}) := P(Y = 1 | X = \mathbf{x})$.

Given (\mathbf{x}_i, y_i) for $i = 1, \dots, n$, assume independence:

$$P(Y_1 = y_1, ..., Y_n = y_n | X_1 = \mathbf{x}_1, ..., X_n = \mathbf{x}_n) = \prod_{i=1}^n P(Y_i = y_i | X_i = \mathbf{x}_i)$$
$$= \prod_{i=1}^n [p(\mathbf{x}_i; \mathbf{w})]^{y_i} [1 - p(\mathbf{x}_i; \mathbf{w})]^{1-y_i}$$

Definition: Maximum Likelihood Estimation (MLE)

$$\max_{\mathbf{w}} \prod_{i=1}^{n} [p(\mathbf{x}_i; \mathbf{w})]^{y_i} [1 - p(\mathbf{x}_i; \mathbf{w})]^{1-y_i}$$

or we can maximize the log-likelihood or equivalently, minimizing the negative log-likelihood

$$\min_{\mathbf{w}} \sum_{i=1}^{n} [-y_i \log p(\mathbf{x}_i; \mathbf{w}) - (1 - y_i) \log(1 - p(\mathbf{x}_i; \mathbf{w}))]$$

Definition: Odds Ratio

Defined as

probability of event probability of non-event

Definition: Logit Transform

The log odds ratio. Assume it is a linear function:

$$\langle \mathbf{x}, \mathbf{w} \rangle = \log \frac{p(\mathbf{x}; \mathbf{w})}{1 - p(\mathbf{x}; \mathbf{w})}$$

where \mathbf{x} is padded with 1 and \mathbf{w} is padded with b.

Definition: Sigmoid Transform

$$p(\mathbf{x}; \mathbf{w}) = \frac{1}{1 + e^{-\langle \mathbf{x}, \mathbf{w} \rangle}}$$

The logit and sigmoid transformation are equivalent.

If we plug in the sigmoid parameterization $p(\mathbf{x}; \mathbf{w}) = \frac{1}{1 + e^{-\langle \mathbf{x}, \mathbf{w} \rangle}}$:

$$\min_{\mathbf{w}} \sum_{i=1}^{n} \log[1 + e^{-\langle \mathbf{x}_i, \mathbf{w} \rangle}] + (1 - y_i) \langle \mathbf{x}_i, \mathbf{w} \rangle$$

Definition: Logistic Loss

If $y_i \in \{\pm 1\}$ in the above, then

$$\min_{\mathbf{w}} \sum_{i=1}^{n} \log[1 + e^{-y_i \langle \mathbf{x}_i, \mathbf{w} \rangle}]$$

Solving Logistic Regression:

$$\min_{\mathbf{w}} \sum_{i=1}^{n} \log[1 + e^{-y_i \langle \mathbf{x}_i, \mathbf{w} \rangle}]$$

We can use gradient descent algorithm: $\mathbf{w} \leftarrow \mathbf{w} - \eta \cdot \nabla_{\mathbf{w}} Loss(\mathbf{w})$.

Prediction:

$$p(\mathbf{x}; \mathbf{w}) = \operatorname{sigmoid}(\langle \mathbf{x}, \mathbf{w} \rangle) = \frac{1}{1 + e^{-\langle \mathbf{x}, \mathbf{w} \rangle}}$$

 $\hat{y} = 1$ if and only if $p(\mathbf{x}; \mathbf{w}) = P(Y = 1 | X = \mathbf{x}) > \frac{1}{2}$ if and only if $\langle \mathbf{x}; \mathbf{w} \rangle > 0$.

Therefore, $\hat{y} = \arg \max_{k} P(Y = k | X = \mathbf{x}).$

The decision boundary remains to be $H := \{ \mathbf{x} : \langle \mathbf{x}, \mathbf{w} \rangle = 0 \}$. This allows us to now predict $\hat{y} = \text{sign}(\langle \mathbf{x}, \mathbf{w} \rangle)$ as before, but with confidence $p(\mathbf{x}; \mathbf{w})$.

Logistic regression estimates the posterior probability $\eta(\mathbf{x}) := P(Y = 1 | X = \mathbf{x})$ under linear log-odds ratio assumption.

This confidence is meaningless if assumption does not hold true.

Classification requires comparing $\eta(\mathbf{x})$ with $\frac{1}{2}$.

4.1 Multiclass Classification

We have class $y \in \{1, ..., c\}$ and we want to learn $\{\mathbf{w}_1, ..., \mathbf{w}_c\}$ for each class.

Definition: Softmax Function

$$P(Y = k | X = \mathbf{x}; \mathbf{W} = [\mathbf{w}_1, \dots, \mathbf{w}_c]) = \frac{e^{\langle \mathbf{x}, \mathbf{w}_k \rangle}}{\sum_{l=1}^c e^{\langle \mathbf{x}, \mathbf{w}_l \rangle}}$$

Training: Use maximum likelihood estimation

$$\min_{\mathbf{w}} \hat{\mathbb{E}} \left[-\log \frac{e^{\langle X, \mathbf{w}_Y \rangle}}{\sum_{l=1}^c e^{\langle X, \mathbf{w}_l \rangle}} \right]$$

Prediction: $\hat{y} = \arg\max_{k} P(Y = k | X = \mathbf{x}; \mathbf{W} = [\mathbf{w}_1, \dots, \mathbf{w}_c])$

Hard-Margin Support Vector Machines

5.1 Margin

Let $H := \{ \mathbf{x} : \langle \mathbf{x}, \mathbf{w} \rangle + b = 0 \}$. What is the distance from a point \mathbf{x}_i to H?

 \mathbf{w} is orthogonal to H. Let \mathbf{x} be any vector in H. The distance = length of the projection of $\mathbf{x}_i - \mathbf{x}$ onto \mathbf{w} .

$$dist(\mathbf{x}_i, H) = \frac{|\langle \mathbf{x}_i - \mathbf{x}, \mathbf{w} \rangle|}{\|\mathbf{w}\|_2} = \frac{|\langle \mathbf{x}_i, \mathbf{w} \rangle - \langle \mathbf{x}, \mathbf{w} \rangle|}{\|\mathbf{w}\|_2} = \frac{|\langle \mathbf{x}_i, \mathbf{w} \rangle + b|}{\|\mathbf{w}\|_2} = \frac{y_i \hat{y}_i}{\|\mathbf{w}\|_2}$$

where the second last equality is from $\mathbf{x} \in H$ and the last equality is from $y_i \hat{y}_i > 0$.

Definition: Margin

The smallest distance to a separating hyperplane H among all separable training data.

$$\min_{i} \frac{y_i \hat{y}_i}{\|\mathbf{w}\|_2} = \min_{i} \frac{|\langle \mathbf{x}_i, \mathbf{w} \rangle + b|}{\|\mathbf{w}\|_2}$$

The margin with respect to a separating hyperplane is the minimum distance to every point.

Our goal is to maximize the margin among all hyperplanes.

$$\max_{\mathbf{w},b} \min_{i} \frac{y_i \hat{y}_i}{\|\mathbf{w}\|_2}, \text{ s.t.} y_i \hat{y}_i > 0, \forall i$$

5.2 Transform to Standard Form

Our current optimization problem is

$$\max_{\mathbf{w}, b: \forall i, y_i \hat{y}_i > 0} \min_{i} \frac{y_i \hat{y}_i}{\|\mathbf{w}\|_2}$$

Both the numerator and denominator are homogeneous in (\mathbf{w}, b) meaning that (\mathbf{w}, b) and $(c\mathbf{w}, cb)$ will have the same loss for c > 0. Varying c will not break the condition $y_i\hat{y}_i > 0$ and can change the numerator arbitrarily.

Let us fix the numerator arbitrarily to 1:

$$\max_{\mathbf{w},b} \frac{1}{\|\mathbf{w}\|_2} \text{ s.t. } \min_i y_i \hat{y}_i = 1$$

Change the max to min and squaring:

$$\min_{\mathbf{w},b} \frac{1}{2} \|\mathbf{w}\|_{2}^{2}$$
s.t. $y_{i}(\langle \mathbf{x}_{i}, \mathbf{w} \rangle + b) \geq 1, \forall i$

5.3 Comparison to Perceptron

Hard-Margin SVM: quadratic programming, unique solution, margin maximizing.

$$\min_{\mathbf{w},b} \quad \frac{1}{2} \|\mathbf{w}\|_{2}^{2}
\text{s.t.} \quad y_{i} \hat{y}_{i} \ge 1, \forall i$$

Perceptron: linear programming, infinitely many solutions, convergence rate depends on maximum margin.

$$\min_{\mathbf{w},b} \quad 0$$
s.t. $y_i \hat{y}_i \ge 1, \forall i$

There are 3 parallel hyperplanes:

$$H := \{\mathbf{x} : \langle \mathbf{x}, \mathbf{w} \rangle + b = 0\}$$

$$H_{+} := \{\mathbf{x} : \langle \mathbf{x}, \mathbf{w} \rangle + b = 1\}$$

$$H_{-} := \{\mathbf{x} : \langle \mathbf{x}, \mathbf{w} \rangle + b = -1\}$$

Support Vectors: Points lie on the supporting hyperplanes H_+ and H_- . If we delete data points not on the H_+ or H_- , H is still uniquely defined.

5.4 Lagrangian Dual

Introducing the Lagrangian multipliers (dual variables) $\alpha \in \mathbb{R}^n$:

$$\min_{\mathbf{w},b} \max_{\alpha \geq 0} \frac{1}{2} \|\mathbf{w}\|_{2}^{2} - \sum_{i} \alpha_{i} [y_{i}(\langle \mathbf{x}_{i}, \mathbf{w} \rangle + b) - 1]$$

$$= \min_{\mathbf{w},b} \begin{cases} +\infty & \text{if } \exists i, y_{i}(\langle \mathbf{x}_{i}, \mathbf{w} \rangle + b) < 1 \text{ (set } \alpha_{i} \text{ as } +\infty) \\ \frac{1}{2} \|\mathbf{w}\|_{2}^{2} & \text{if } \forall i, y_{i}(\langle \mathbf{x}_{i}, \mathbf{w} \rangle + b) \geq 1 \text{ (set } \alpha_{i} \text{ as } 0 \text{ for all } i) \end{cases}$$

$$= \min_{\mathbf{w},b} \frac{1}{2} \|\mathbf{w}\|_{2}^{2}, \text{ s.t. } y_{i}(\langle \mathbf{x}_{i}, \mathbf{w} \rangle + b) \geq 1, \forall i$$

Transfer a constrained optimization problem to an unconstrained one on \mathbf{w}, b .

In the above, we proved the hard-margin SVM is equivalent to

$$\min_{\mathbf{w},b} \max_{\alpha \ge \mathbf{0}} \frac{1}{2} \|\mathbf{w}\|_2^2 - \sum_i \alpha_i [y_i(\langle \mathbf{x}_i, \mathbf{w} \rangle + b) - 1]$$

Swap min with max (when primal is convex (strong duality)):

$$\max_{\alpha \ge \mathbf{0}} \min_{\mathbf{w}, b} \frac{1}{2} \|\mathbf{w}\|_2^2 - \sum_i \alpha_i [y_i(\langle \mathbf{x}_i, \mathbf{w} \rangle + b) - 1]$$

Solve the inner unconstrained problem by setting the derivative to 0:

$$\frac{\partial}{\partial \mathbf{w}} = \mathbf{w} - \sum_{i} \alpha_{i} y_{i} \mathbf{x}_{i} = 0, \ \frac{\partial}{\partial b} = -\sum_{i} \alpha_{i} y_{i} = 0$$

Plug w in the loss to eliminate inner problem:

$$Loss(\alpha) = \min_{\mathbf{w}, b} \frac{1}{2} \|\mathbf{w}\|_{2}^{2} - \sum_{i} \alpha_{i} [y_{i}(\langle \mathbf{x}_{i}, \mathbf{w} \rangle + b) - 1]$$
$$= \frac{1}{2} \left\| \sum_{i} \alpha_{i} y_{i} \mathbf{x}_{i} \right\|_{2}^{2} - \left\langle \sum_{i} \alpha_{i} y_{i} \mathbf{x}_{i}, \sum_{i} \alpha_{i} y_{i} \mathbf{x}_{i} \right\rangle - b \sum_{i} \alpha_{i} y_{i} + \sum_{i} \alpha_{i}$$

That is, $\max_{\alpha \geq \mathbf{0}} \sum_{i} \alpha_{i} - \frac{1}{2} \left\| \sum_{i} \alpha_{i} y_{i} \mathbf{x}_{i} \right\|_{2}^{2}$ such that $\sum_{i} \alpha_{i} y_{i} = 0$. Changing it to a minimization problem and expanding the norm:

$$\min_{\alpha \ge \mathbf{0}} - \sum_{i} \alpha_i + \frac{1}{2} \sum_{i} \sum_{j} \alpha_i \alpha_j y_i y_j \left\langle \mathbf{x}_i, \mathbf{x}_j \right\rangle, \text{ s.t. } \sum_{i} \alpha_i y_i = 0$$

Comparing the Primal to Dual: Primal variables: $\mathbf{w} \in \mathbb{R}^d, b \in \mathbb{R}$, n primal inequalities, and 0 primal equalities.

Dual variables: $\alpha \in \mathbb{R}^n$ where each α_i corresponds to a data sample, n dual inequalities, and 1 dual equality.

Soft-Margin Support Vector Machine

Suppose the data is not linearly separable, we want to penalize $y(\langle \mathbf{x}, \mathbf{w} \rangle + b) < 1$, i.e. small or negative $y\hat{y}$.

Recall that maximizing the margin is given by:

$$\max_{\mathbf{w}} \frac{1}{\|w\|_2}$$

or we can equivalently minimize

$$\min_{\mathbf{w}} \frac{1}{2} \left\| w \right\|_2^2$$

where the $\frac{1}{2}$ is for ease of gradient calculation.

6.1 Hinge Loss

Definition: Hinge Loss

Let $y \in \{\pm 1\}$ be the true label and $\hat{y} \in \mathbb{R}$ be the predicted confidence, then

$$l_{hinge}(y\hat{y}) = (1 - y\hat{y})^{+} = \begin{cases} 1 - y\hat{y} & \text{if } y\hat{y} < 1\\ 0 & \text{otherwise} \end{cases}$$

We want to penalize two types of points: one that is misclassified and one that is correctly classified, but too close to the decision boundary (within margin). Balance between margin maximization and the hinge loss.

Definition: Soft Margin SVM

$$\min_{\mathbf{w},b} \quad \frac{1}{2} \|\mathbf{w}\|_{2}^{2} + C \cdot \sum_{i=1}^{n} (1 - y_{i} \hat{y}_{i})^{+}$$
s.t.
$$\hat{y}_{i} = \langle \mathbf{x}_{i}, \mathbf{w} \rangle + b, \forall i$$

 $\frac{1}{2} \|w\|_2^2$ is the margin maximization, $(1 - y_i \hat{y}_i)^+$ is the error penalty, and C is a hyperparameter to control tradeoff.

Hard-margin SVM has hard constraints which we must respect while the soft-margin SVM has a soft penalty: the more you deviate, the heavier the penalty.

Our goal: minimize over \mathbf{w}, b . Recall for any event A, then $P[A] = \mathbb{E}\mathbb{I}[A]$ where $\mathbb{I}[A] = 1$ if A happens, and 0 otherwise.

$$P(Y \neq \operatorname{sign}(\hat{Y})) = P(Y\hat{Y} \leq 0) = \mathbb{E}[\mathbb{I}[Y\hat{Y} \leq 0]] := \mathbb{E}[l_{0-1}(Y\hat{Y})]$$

So, our goal is to minimize $\mathbb{E}[l_{0-1}(Y\hat{Y})]$. Even with linear predictors, minimizing the above 0-1 error is NP-hard.

- The loss is not continuous at 0.
- The gradient of the loss is 0 almost surely.

Therefore, we need to consider a surrogate loss, e.g., the hinge loss.

Definition: Bayes' Rule

Given an instance \mathbf{x} , Bayes rule is given by

$$\eta(\mathbf{x}) := \underset{\hat{y} \in \mathbb{R}}{\arg\min} \, \mathbb{E}[l_{0-1}(Y\hat{y})|X = \mathbf{x}]$$

$$\eta(\mathbf{x}) = \underset{\hat{y} \in \mathbb{R}}{\arg\min} \, \mathbb{E}[l_{0-1}(Y\hat{y})|X = \mathbf{x}]$$
$$= \underset{\hat{y} \in \mathbb{R}}{\arg\min} \, P[Y \neq \text{sign}(\hat{y})|X = \mathbf{x}]$$

Definition: Classification Calibrated

We say a loss $l(y\hat{y})$ is classification calibrated if and only if for all \mathbf{x} ,

$$\hat{y}(\mathbf{x}) := \underset{\hat{y} \in \mathbb{R}}{\operatorname{arg\,min}} \, \mathbb{E}[l(Y\hat{y})|X = \mathbf{x}]$$

has the same sign as the Bayes rule $\eta(\mathbf{x}) = \arg\min_{\hat{y} \in \mathbb{R}} \mathbb{E}[l_{0-1}(Y\hat{y})|X = \mathbf{x}].$

Theorem (Characterization Under Convexity)

Any convex loss l is classification calibrated if and only if l is differentiable at 0 and l'(0) < 0.

Therefore, the classifier that minimizes the expected hinge loss minimizes the expected zero-one loss.

6.2 Lagrangian Dual

Recall the soft-margin SVM: $\min_{\mathbf{w},b} \frac{1}{2} \|\mathbf{w}\|_2^2 + C \cdot \sum_{i=1}^n (1 - y_i(\langle \mathbf{x}_i, \mathbf{w} \rangle + b))^+$.

Apply $C \cdot (t)^+ := \max\{Ct, 0\} = \max_{0 \le \alpha \le C} \alpha t$:

$$\min_{\mathbf{w},b} \max_{0 \le \alpha \le C} \frac{1}{2} \|\mathbf{w}\|_{2}^{2} + C \cdot \sum_{i=1}^{n} \alpha_{i} [1 - y_{i}(\langle \mathbf{x}_{i}, \mathbf{w} \rangle + b)]$$

We can swap min with max because strong duality holds due to convexity:

$$\max_{0 \le \alpha \le C} \min_{\mathbf{w}, b} \frac{1}{2} \|\mathbf{w}\|_{2}^{2} + C \cdot \sum_{i=1}^{n} \alpha_{i} [1 - y_{i}(\langle \mathbf{x}_{i}, \mathbf{w} \rangle + b)]$$

Solving the inner unconstrained problem by setting gradient to 0:

$$\frac{\partial}{\partial \mathbf{w}} = \mathbf{w} - \sum_{i=1}^{n} \alpha_i y_i \mathbf{x}_i = \mathbf{0}, \frac{\partial}{\partial b} = -\sum_{i=1}^{n} \alpha_i y_i = 0$$

We can plug back to eliminate the inner problem to get the dual:

$$\max_{0 \le \alpha \le C} \sum_{i} \alpha_{i} - \frac{1}{2} \left\| \sum_{i} \alpha_{i} y_{i} \mathbf{x}_{i} \right\|_{2}^{2} \text{ s.t. } \sum_{i} \alpha_{i} y_{i} = 0$$

We can change max to min and expand the norm:

$$\min_{0 \le \alpha \le C} \frac{1}{2} \sum_{i} \sum_{j} \alpha_i \alpha_j y_i y_j \langle \mathbf{x}_i, \mathbf{x}_j \rangle - \sum_{i} \alpha_i \text{ s.t. } \sum_{i} \alpha_i y_i = 0$$

If $C \to \infty$, then it is a hard-margin SVM. If $C \to 0$, then it is a constant classifier.

6.3 Complementary Slackness

The following relation is to introduce dual variables:

$$C \cdot (t)^+ = \max\{Ct, 0\} = \max_{0 \le \alpha \le C} \alpha t = \alpha^* t$$

The complementary slackness conditions are

•
$$t > 0 \implies \alpha^* = C, \alpha^* = C \implies t > 0$$

•
$$t < 0 \implies \alpha^* = 0, \alpha^* = 0 \implies t \le 0$$

6.4 Recovering Parameters

We can use α to determine the location of data points.

Recovering $\mathbf{w}^* := \sum_i \alpha_i^* y_i \mathbf{x}_i$.

We normally set C large enough so that there is at least one data point on $H_{\pm 1}$ or $y_i\hat{y}_i=1$.

This point can be used to recover b^* : $y(\langle \mathbf{x}, \mathbf{w}^* \rangle + b^*) = 1 \implies b^* = y - \langle \mathbf{x}, \mathbf{w}^* \rangle$.

6.5 Training by Gradient Descent

$$\min_{\mathbf{w},b} \frac{1}{2\lambda} \|\mathbf{w}\|_{2}^{2} + \frac{1}{n} \sum_{i=1}^{n} l(y_{i} \hat{y}_{i})$$

Gradient descent:

$$\mathbf{w} \leftarrow \mathbf{w} - \eta \left[\frac{\mathbf{w}}{\lambda} + \frac{1}{n} \sum_{i=1}^{n} l'(y_i \hat{y}_i) y_i \mathbf{x}_i \right]$$

The gradient of the hinge loss is

$$l'_{hinge}(t) = \begin{cases} -1 & t < 1 \\ 0 & t > 1 \text{ while we choose } l'_{perceptron}(t) = \begin{cases} -1 & t \le 0 \\ 0 & t > 0 \end{cases}$$
$$[-1, 0] \quad t = 1$$

Reproducing Kernels

Linear classifiers may not be able to separate the data, such as the XOR dataset. We need a more complex classifier.

Definition: Quadratic Classifier

$$f(\mathbf{x}) = \langle \mathbf{x}, Q\mathbf{x} \rangle + \sqrt{2} \langle \mathbf{x}, \mathbf{p} \rangle + b$$

We predict as before: $\hat{y} = \text{sign}(f(\mathbf{x}))$ where $\mathbf{x} \in \mathbb{R}^d$. The weights can be learned: $Q \in \mathbb{R}^{d \times d}$, $\mathbf{p} \in \mathbb{R}^d$, $b \in R$.

Note that setting $Q = \mathbf{0}$ reduces to the linear case.

Definition: Vectorization

For any matrix $A \in \mathbb{R}^{m \times n}$, let \vec{A} be the vectorization operation from $\mathbb{R}^{m \times n} \to \mathbb{R}^{mn}$ by copying each row and appending the bottom as a column vector.

$$f(\mathbf{x}) = \langle \mathbf{x}, Q\mathbf{x} \rangle + \sqrt{2} \langle \mathbf{x}, \mathbf{p} \rangle + b$$
$$= \langle \mathbf{x}\mathbf{x}^T, Q \rangle + \langle \sqrt{2}\mathbf{x}, \mathbf{p} \rangle + b$$
$$= \langle \phi(\mathbf{x}), \mathbf{w} \rangle$$

We define the feature map

$$\phi(\mathbf{x}) = \begin{bmatrix} \vec{\mathbf{x}} \vec{\mathbf{x}}^T \\ \sqrt{2}\mathbf{x}, 1 \end{bmatrix}$$

where $\mathbf{x} \in \mathbb{R}^d \mapsto \phi(\mathbf{x}) \in \mathbb{R}^{d \times d + d + 1}$. The weights to be learned:

$$\mathbf{w} = \begin{bmatrix} \vec{Q} \\ \mathbf{p} \\ b \end{bmatrix} \in \mathbb{R}^{d \times d + d + 1}$$

The inner product of two matrices $\langle A, B \rangle$ is the inner product of the vectorization of A and B.

The function is nonlinear in \mathbf{x} , but linear in $\phi(\mathbf{x})$.

7.1 The Kernel Trick

The feature map $\phi: \mathbb{R}^d \to \mathbb{R}^{d \times d + d + 1}$ has huge dimensions. We do not have to operate in this high-dimensional feature space explicitly. In the dual form of the SVM, all we need is the inner product.

$$\langle \phi(\mathbf{x}), \phi(\mathbf{z}) \rangle = \left\langle \begin{bmatrix} \vec{\mathbf{x}} \vec{\mathbf{x}}^T \\ \sqrt{2} \mathbf{x} \\ 1 \end{bmatrix}, \begin{bmatrix} \vec{\mathbf{z}} \vec{\mathbf{z}}^T \\ \sqrt{2} \mathbf{z} \\ 1 \end{bmatrix} \right\rangle = (\langle \mathbf{x}, \mathbf{z} \rangle)^2 + 2 \langle \mathbf{x}, \mathbf{z} \rangle + 1 = (\langle \mathbf{x}, \mathbf{z} \rangle + 1)^2$$

We see that the inner product can be computed by the original vectors and can be done in O(d) runtime.

We saw that given a feature map $\phi : \mathcal{X} \to \mathcal{H}$, $\langle \phi(\mathbf{x}), \phi(\mathbf{z}) \rangle = k(\mathbf{x}, \mathbf{z})$. Can we do the converse: given $k : \mathcal{X} \times \mathcal{X} \to \mathbb{R}$, does there exist $\phi : \mathcal{X} \to \mathcal{H}$.

Definition: (Reproducing) Kernel

We call $k: \mathcal{X} \times \mathcal{X} \to \mathbb{R}$ a (reproducing) kernel if and only if there exists some feature transform $\phi: \mathcal{X} \to \mathcal{H}$ so that $\langle \phi(\mathbf{x}), \phi(\mathbf{z}) \rangle = k(\mathbf{x}, \mathbf{z})$.

Choosing a feature transform ϕ determines the unique corresponding kernel k and choosing a kernel k determines some feature transform ϕ .

Definition: Positive Semi-Definite (PSD)

For any $\alpha \in \mathbb{R}^n$,

$$\langle \alpha, K\alpha \rangle = \sum_{i=1}^{n} \sum_{j=1}^{n} \alpha_i \alpha_j K_{ij} \ge 0$$

Theorem (Mercer)

 $k: \mathcal{X} \times \mathcal{X} \to \mathbb{R}$ is a kernel if and only if for any $n \in \mathbb{N}$, for any $\mathbf{x}_1, \dots, \mathbf{x}_n \in \mathcal{X}$, the kernel matrix K such that $K_{ij} := k(\mathbf{x}_i, \mathbf{x}_j)$ is symmetric and positive semi-definite.

Therefore, we do not need to find ϕ , the kernel trick allows us to just use a kernel.

Definition: Polynomial Kernel

Let p be a hyper-parameter, then

$$k(\mathbf{x}, \mathbf{z}) = (\langle \mathbf{x}, \mathbf{z} \rangle + 1)^p$$

where a larger p means a higher-degree polynomial mapping ϕ .

Definition: Gaussian Kernel

$$k(\mathbf{x}, \mathbf{z}) = \exp(-\|\mathbf{x} - \mathbf{z}\|_2^2 / \sigma)$$

Definition: Laplace Kernel

Let σ be a hyper-parameter, then

$$k(\mathbf{x}, \mathbf{z}) = \exp(-\|\mathbf{x} - \mathbf{z}\|_2 / \sigma)$$

where the larger σ means a smoother ϕ , i.e. $\phi(x_1)$ and $\phi(x_2)$ will not differ too much for close x_1, x_2 .

Proposition

If k is a kernel, then λk is a kernel for any $\lambda \geq 0$.

Proposition

If k_1 and k_2 are kernels, then $k_1 + k_2$ is a kernel.

Proposition

If k_1 and k_2 are kernels, then k_1k_2 is a kernel.

Proposition

If k_t are kernels, then the limit $\lim_t k_t$, when it exists, is a kernel.

7.2 Kernel SVM

The dual for the soft-margin SVM is now

$$\min_{C \ge \alpha \ge 0} -\sum_{i} \alpha_{i} + \frac{1}{2} \sum_{i} \sum_{j} \alpha_{i} \alpha_{j} y_{i} y_{j} k(\mathbf{x}_{i}, \mathbf{x}_{j})$$
s.t.
$$\sum_{i} \alpha_{i} y_{i} = 0$$

We can solve for $\alpha^* \in \mathbb{R}^n$ and recover $\mathbf{w}^* = \sum_{i=1}^n \alpha_i^* y_i \phi(\mathbf{x}_i)$. We do not know ϕ , so for testing, we only need to compute

$$f(\mathbf{x}) = \langle \phi(\mathbf{x}), \mathbf{w}^* \rangle = \left\langle \phi(\mathbf{x}), \sum_{i=1}^n \alpha_i^* y_i \phi(\mathbf{x}_i) \right\rangle = \sum_{i=1}^n \alpha_i^* y_i k(\mathbf{x}, \mathbf{x}_i)$$

We only need to know the dual vector α^* , the training set and kernel k.

Gradient Descent

Many machine learning models can be formulated as an optimization problem.

Consider an unconstrained optimization

$$\min_{x} f(x)$$

Assume that f is differentiable with gradient $\nabla f(x)$. Denote the optimal criterion value by $f^* = \min_x f(x)$ and an optimal solution x^* .

Gradient Descent Template

Choose initial $x^{(0)} \in \mathbb{R}^n$ and repeat

$$x^{(k)} = x^{(k-1)} - \underbrace{t}_{\text{step size}} \cdot \nabla f(x^{(k-1)})$$

for
$$k = 1, 2, ...$$

The **negative gradient** is the steepest decreasing direction at that point. So if the step size is small and the function is convex, the algorithm will reach the minimizer.

8.1 Gradient Descent in ML Models

Perceptron:

$$\min_{\mathbf{w}} -\frac{1}{n} \sum_{i=1}^{n} y_i \langle \mathbf{w}, \mathbf{x}_i \rangle \mathbb{I}[\text{mistake on } \mathbf{x}_i]$$

The gradient descent update is

$$\mathbf{w} \leftarrow \mathbf{w} + t \left[\frac{1}{n} \sum_{i=1}^{n} y_i \mathbf{x}_i \mathbb{I}[\text{mistake on } \mathbf{x}_i] \right]$$

Stochastic gradient descent update is

$$\mathbf{w} \leftarrow \mathbf{w} + ty_I \mathbf{x}_I \mathbb{I}[\text{mistake on } \mathbf{x}_i]$$

for some random index I.

Soft-Margin SVM:

$$\min_{\mathbf{w},b} \frac{1}{2\lambda} \|\mathbf{w}\|_2^2 + \frac{1}{n} \sum_{i=1}^n l_{hinge}(y_i \hat{y}_i) \text{ s.t. } \hat{y}_i = \langle \mathbf{x}_i, \mathbf{w} \rangle + b$$

The gradient descent update is

$$\mathbf{w} \leftarrow \mathbf{w} - t \left[\frac{\mathbf{w}}{\lambda} + \frac{1}{n} \sum_{i=1}^{n} l'_{hinge}(y_i \hat{y}_i) y_i \mathbf{x}_i \right]$$
$$b \leftarrow b - t \left[\frac{1}{n} \sum_{i=1}^{n} l'_{hinge}(y_i \hat{y}_i) y_i \right]$$

8.2 Convergence Analysis

8.2.1 Taylor Series

Consider the min operation on the Taylor expansion of f locally at x, where x is the current iterate:

$$\min_{y} f(y) \approx \min_{y} \left[f(x) + \nabla f(x)^{T} (y - x) + \frac{1}{2t} \|y - x\|_{2}^{2} \right]$$

Taking the gradient with respect to y,

$$\frac{\partial f(y)}{\partial y} = 0 \implies \nabla f(x) + \frac{1}{t}(y - x) = 0 \implies y^* = x - t \cdot \nabla f(x)$$

Geometrically, Choose the next point $y = x^+$ to minimize the RHS:

$$x^{(i+1)} = \operatorname*{arg\,min}_{y} f(x^{(i)}) + \nabla f(x^{(i)})^{T} (y - x) + \frac{1}{2t} \left\| y - x^{(i)} \right\|_{2}^{2}$$

This will diverge if t is too large. If t is too small, it will be too slow.

8.2.2 Convex Case

Definition: Convex Function (Differentiable)

For any $x, y \in \mathbb{R}^n$,

$$f(y) \ge f(x) + \nabla f(x)^T (y - x)$$

For thin convex functions, L is bigger and choosing a step size that is too large will jump to the other side of the function, so it will diverge. For wider convex functions, L can be smaller and we are able to choose a larger step size.

Theorem (Convergence Rate For Convex Case)

Assume f is convex and differentiable, with $dom(f) = \mathbb{R}^n$, and ∇f is L-Lipschitz continuous: $L\mathbf{I} - \nabla^2 f(x)$ is positive semi-definite $(L\mathbf{I} \succeq \nabla^2 f(x))$. Gradient descent with fixed step size $t \leq \frac{1}{L}$ satisfies

$$f(x^{(k)}) - f^* \le \frac{\left\|x^{(0)} - x^*\right\|_2^2}{2tk}$$

We say gradient descent has convergence rate O(1/k). That is, it finds ϵ -suboptimal point in $O(1/\epsilon)$ iterations.

Proof. For any y, perform a quadratic expansion and obtain:

$$f(y) \leq f(x) + \nabla f(x)^{T} (y - x) + \frac{1}{2} (y - x)^{T} \nabla^{2} f(x) (y - x)$$

$$\leq f(x) + \nabla f(x)^{T} (y - x) + \frac{1}{2} L \|y - x\|_{2}^{2} \qquad (L\mathbf{I} \succeq \nabla^{2} f(x))$$

$$f(x^{+}) \leq f(x) + \nabla f(x)^{T} (x^{+} - x) + \frac{1}{2} L \|x^{+} - x\|_{2}^{2} \qquad (y = x^{+} = x - t \nabla f(x))$$

$$= f(x) + \nabla f(x)^{T} (x - t \nabla f(x) - x) + \frac{1}{2} L \|x - t \nabla f(x) - x\|_{2}^{2}$$

$$= f(x) - \left(1 - \frac{1}{2} Lt\right) t \|\nabla f(x)\|_{2}^{2} \qquad (1)$$

$$\leq f(x) - \frac{1}{2} t \|\nabla f(x)\|_{2}^{2} \qquad (t \leq 1/L)$$

Each update decreases the function value by at least $\frac{1}{2}t \|\nabla f(x)\|_2^2$. Since f is convex

$$f(x^*) \ge f(x) + \nabla f(x)^T (x^* - x) \implies f(x) \le f(x^*) + \nabla f(x)^T (x - x^*)$$

Plug in (1)

$$f(x^{+}) \leq f(x^{*}) + \nabla f(x)^{T}(x - x^{*}) - \frac{t}{2} \|\nabla f(x)\|_{2}^{2}$$

$$f(x^{+}) - f(x^{*}) \leq \frac{1}{2t} (2t\nabla f(x)^{T}(x - x^{*}) - t^{2} \|\nabla f(x)\|_{2}^{2})$$

$$f(x^{+}) - f(x^{*}) \leq \frac{1}{2t} (2t\nabla f(x)^{T}(x - x^{*}) - t^{2} \|\nabla f(x)\|_{2}^{2} - \|x - x^{*}\|_{2}^{2} + \|x - x^{*}\|_{2}^{2})$$

$$f(x^{+}) - f(x^{*}) \leq \frac{1}{2t} (\|x - x^{*}\|_{2}^{2} - \|x - t\nabla f(x) - x^{*}\|_{2}^{2})$$

$$f(x^{+}) - f(x^{*}) \leq \frac{1}{2t} (\|x - x^{*}\|_{2}^{2} - \|x^{+} - x^{*}\|_{2}^{2})$$

Sum over all iterations:

$$\sum_{i=1}^{k} (f(x^{(i)}) - f(x^*)) \le \sum_{i=1}^{k} \frac{1}{2t} (\|x^{(i-1)} - x^*\|_2^2 - \|x^{(i)} - x^*\|_2^2)$$

$$= \frac{1}{2t} (\|x^{(0)} - x^*\|_2^2 - \|x^{(k)} - x^*\|_2^2)$$

$$\le \frac{1}{2t} \|x^{(0)} - x^*\|_2^2$$

which implies

$$f(x^{(k)}) \le \frac{1}{k} \sum_{i=1}^{k} f(x^{(i)}) \le f(x^*) + \frac{\left\|x^{(0)} - x^*\right\|_2^2}{2tk}$$

The first inequality holds because $f(x^{(i)})$ decreases as i increases.

8.2.3 Strong Convex Case

Definition: m-Strong Convexity

A function f is m-strong convex if $f(x) - m ||x||_2^2$ is convex: $L\mathbf{I} \succeq \nabla^2 f(x) \succeq m\mathbf{I}$.

Theorem (Convergence Rate For Strong Convexity)

Let f be m-strongly convex and L-smooth. Gradient descent with fixed step size $t \leq \frac{2}{m+L}$ satisfies

$$f(x^{(k)}) - f^* \le \gamma^k \frac{L}{2} \|x^{(0)} - x^*\|_2^2$$

where $0 \le \gamma \le 1$.

The rate under strong convexity is $O(\gamma^k)$, which is exponentially fast. That is, it finds ϵ -suboptimal point, i.e. $f(x^{(k)}) - f^* < \epsilon$, in $O(\log(1/\epsilon))$ iterations.

Asking for optimality is too much. Focus on $\|\nabla f(x)\|_2 \leq \epsilon$.

8.2.4 Nonconvex Case

Assume f is differentiable with L-Lipschitz gradient and f is nonconvex.

Theorem (Convergence Rate For Nonconvex Case)

Gradient descent with fixed step size $t \leq \frac{1}{L}$ satisfies

$$\min_{i=0,\dots,k} \left\| \nabla f(x^{(i)}) \right\|_2^2 \le \sqrt{\frac{2(f(x^{(0)}) - f^*)}{t(k+1)}}$$

The gradient descent rate is $O(1/\sqrt{k})$ or $O(1/\epsilon^2)$, even in the nonconvex case for finding stationary points.

This rate cannot be improved by any deterministic algorithm.

8.3 Stochastic Gradient Descent

Consider the decomposable optimization (n is very large)

$$\min_{w} \frac{1}{n} \sum_{i=1}^{n} f_i(w)$$

Gradient descent has $w^+ = w - w - t \frac{1}{n} \sum_{i=1}^n \nabla f_i(x)$, but stochastic gradient descent has

$$w^+ = w - t \cdot \nabla f_I(x)$$

where I is a random index. The expectation of gradient descent is the same as stochastic gradient descent.

For gradient descent, the step size $t \leq 1/L$ and the convergence rate is $O\left(\frac{n}{\epsilon}\right)$. For stochastic gradient descent, the step size is t=1/k for $k=1,2,\ldots$ and the convergence rate is $O\left(\frac{1}{\sqrt{k}}\right)$. The randomness leads to a large variance of estimation of the gradient. Thus, SGD requires more iterations.

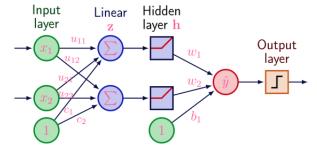
Part II Neural Networks

Multilayer Perceptron/Fully Connected Neural Networks

9.1 Introduction

We have seen that no linear classifier can separate the XOR dataset, so it underfits the data. We can fix the input representation using a richer model (e.g. quadratic classifier) or keep it linear, but use a richer input representation. The two approaches are equivalent for certain classifiers, using a reproducing kernel.

A neural network can learn the feature map simultaneously with the linear classifier.



First linear transformation: $\mathbf{z} = U\mathbf{x} + \mathbf{c}$ where $U \in \mathbb{R}^{2 \times 2}, c \in \mathbb{R}^2$.

Element-wise nonlinear activation function: $\mathbf{h} = \sigma(\mathbf{z})$.

Second linear transformation: $\hat{y} = \langle \mathbf{h}, \mathbf{w} \rangle + b$.

Output layer: $\operatorname{sign}(\hat{y})$ or $\operatorname{sigmoid}(\hat{y})$.

XOR Dataset: Consider a 2-layer NN:

- 1st linear transformation: $\mathbf{z} = U\mathbf{x} + \mathbf{c}$ where $U \in \mathbb{R}^{2\times 2}, \mathbf{c} \in \mathbb{R}^2$
- Element-wise nonlinear activation: $\mathbf{h} = \sigma(\mathbf{z})$; use the Rectified Linear Unit (ReLU): $\sigma(t) = t_+ := \max\{t, 0\}$

• 2nd linear transformation: $\hat{y} = \langle \mathbf{h}, \mathbf{w} \rangle + b$

$$U = \begin{bmatrix} 1 & 1 \\ 1 & 1 \end{bmatrix}, \mathbf{c} = \begin{bmatrix} 0 \\ -1 \end{bmatrix}, \mathbf{w} = \begin{bmatrix} 2 \\ -4 \end{bmatrix}, b = -1$$

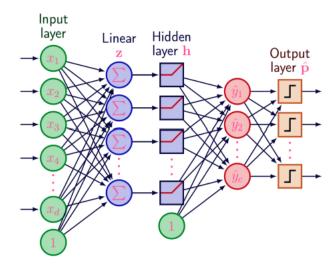
•
$$\mathbf{x}_1 = (0,0), y = - \implies \mathbf{z}_1 = (0,-1), \mathbf{h}_1 = (0,0) \implies \hat{y} = -1 \checkmark$$

•
$$\mathbf{x}_2 = (1,0), y = + \implies \mathbf{z}_2 = (1,0), \mathbf{h}_1 = (1,0) \implies \hat{y} = +1 \checkmark$$

•
$$\mathbf{x}_3 = (0,1), y = + \implies \mathbf{z}_3 = (1,0), \mathbf{h}_1 = (1,0) \implies \hat{y} = +1 \checkmark$$

•
$$\mathbf{x}_4 = (1,1), y = - \implies \mathbf{z}_4 = (2,1), \mathbf{h}_1 = (2,1) \implies \hat{y} = -1 \checkmark$$

9.2 Multiclass Classification



$$\mathbf{z} = U\mathbf{x} + \mathbf{c}, \mathbf{h} = \sigma(\mathbf{z}), \quad \hat{\mathbf{y}} = W\mathbf{h} + \mathbf{b}, \hat{\mathbf{p}} = \operatorname{softmax}(\hat{\mathbf{y}})$$
learning feature \mathbf{h} learning linear classifier by logistic regression

9.3 Activation Functions

Definition: Sigmoid

$$\operatorname{sgm}(t) = \frac{1}{1 + \exp(-t)}$$

Definition: Hyperbolic Tangent

$$\tanh(t) = 1 - 2\operatorname{sgm}(2t)$$

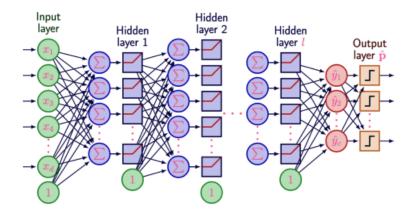
Definition: Rectified Linear Unit (ReLU)

$$ReLU(t) = t_+ := \max\{t, 0\}$$

Definition: Exponential Linear Unit

$$ELU(t) = (t)_{+} + (t)_{-}(\exp(t) - 1)$$

9.4 Training



$$\hat{\mathbf{p}} = f(\mathbf{x}; \mathbf{w})$$

We need a loss function l to measure the difference between the prediction $\hat{\mathbf{p}}$ and the truth y. We can use squared loss $\|\hat{\mathbf{p}} - \mathbf{y}\|_2^2$ for regression or log-loss $-\log \hat{p}_y$ for classification.

We also need a training set $\mathcal{D} = \{(\mathbf{x}_i, y_i) : i = 1, ..., n\}$ to train the weights \mathbf{w} .

Stochastic Gradient Descent: Let $[l \circ f](\mathbf{x}_i, y_i; \mathbf{w}) := l[f(\mathbf{x}_i; \mathbf{w}), y_i]$

$$\min_{\mathbf{w}} \frac{1}{n} \sum_{i=1}^{n} [l \circ f](\mathbf{x}_i, y_i; \mathbf{w})$$

The update is

$$\mathbf{w} \leftarrow \mathbf{w} - \eta \cdot \frac{1}{n} \sum_{i=1}^{n} \nabla[l \circ f](\mathbf{x}_i, y_i; \mathbf{w})$$

However, each iteration requires going through the entire training set, so we can just take a random minibatch $B \subseteq \{1, ..., n\}$.

$$\mathbf{w} \leftarrow \mathbf{w} - \eta \cdot \frac{1}{|B|} \sum_{i \in B} \nabla[l \circ f](\mathbf{x}_i, y_i; \mathbf{w})$$

There is trade-off between variance and computation time using randomness.

The learning rate decay decreases every few epochs:
$$\eta_t = \begin{cases} \eta_0 & t \le t_0 \\ \eta_0/10 & t_0 < t \le t_1 \\ \eta_0/100 & t_1 < t \end{cases}$$

Sublinear decay: $\eta_t = \eta_0/(1+ct)$ or $\eta_t = \eta_0/\sqrt{1+ct}$.

9.4.1 Computing the Gradient

The forward pass of the model is (these are calculated):

$$\mathbf{x} = \text{input}$$

$$\mathbf{z} = W\mathbf{x} + \mathbf{b}_1$$

$$\mathbf{h} = \text{ReLU}(\mathbf{z})$$

$$\boldsymbol{\theta} = U\mathbf{h} + \mathbf{b}_2$$

$$J = \frac{1}{2} \|\boldsymbol{\theta} - \mathbf{y}\|_2^2$$

The parameters to be learned are $W, \mathbf{b}_1, U, \mathbf{b}_2$. The network's gradient are $\frac{\partial J}{\partial W}, \frac{\partial J}{\partial \mathbf{b}_1}, \frac{\partial J}{\partial U}, \frac{\partial J}{\partial \mathbf{b}_2}$. The gradient of ReLU is

$$ReLU'(x) = \begin{cases} 1 & x > 0 \\ 0 & \text{otherwise} \end{cases}$$

Proposition

If $J: \mathbb{R}^{m \times n} \to \mathbb{R}$, then $\frac{\partial J(X)}{\partial X} \in \mathbb{R}^{m \times n}$. The output of the gradient is the same shape as the input.

Definition: Hadamard Product

If $a, b \in \mathbb{R}^n$, then

$$a \odot b = \begin{bmatrix} a_1b_1 \\ a_2b_2 \\ \vdots \\ a_nb_n \end{bmatrix}$$

We can use the chain rule: the backward pass (back-propagation) of the model (⊙ is the

Hadamard product):

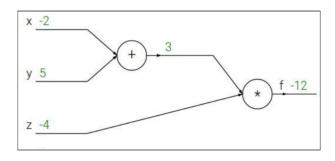
$$\frac{\partial J}{\partial \boldsymbol{\theta}} = \boldsymbol{\theta} - \mathbf{y}
\frac{\partial J}{\partial U} = \frac{\partial J}{\partial \boldsymbol{\theta}} \frac{\partial \boldsymbol{\theta}}{\partial U} = (\boldsymbol{\theta} - \mathbf{y}) \mathbf{h}^{T}$$
(*)
$$\frac{\partial J}{\partial \mathbf{b}_{2}} = \frac{dJ}{d\boldsymbol{\theta}} \frac{\partial \boldsymbol{\theta}}{\partial \mathbf{b}_{2}} = \boldsymbol{\theta} - \mathbf{y}$$
(*)
$$\frac{\partial J}{\partial \mathbf{h}} = \frac{\partial J}{\partial \boldsymbol{\theta}} \frac{\partial \boldsymbol{\theta}}{\partial \mathbf{h}} = U^{T}(\boldsymbol{\theta} - \mathbf{y})$$

$$\frac{dJ}{d\mathbf{z}} = \frac{\partial J}{\partial \mathbf{h}} \frac{\partial \mathbf{h}}{\partial \mathbf{z}} = U^{T}(\boldsymbol{\theta} - \mathbf{y}) \odot \operatorname{ReLU}'(\mathbf{z})$$

$$\frac{\partial J}{\partial W} = \frac{\partial J}{\partial \mathbf{z}} \frac{\partial \mathbf{z}}{\partial W} = (U^{T}(\boldsymbol{\theta} - \mathbf{y}) \odot \operatorname{ReLU}'(\mathbf{z})) \mathbf{x}^{T}$$
(*)
$$\frac{\partial J}{\partial \mathbf{b}_{1}} = \frac{\partial J}{\partial \mathbf{z}} \frac{\partial \mathbf{z}}{\partial \mathbf{b}_{1}} = U^{T}(\boldsymbol{\theta} - \mathbf{y}) \odot \operatorname{ReLU}'(\mathbf{z})$$
(*)

Existing frameworks will memorize the computational graph in the forward process and calculate the back-propagation automatically for you.

Example: f(x, y, z) = (x + y)z with x = -2, y = 5, z = -4.



q = x + y, $\frac{\partial q}{\partial x} = 1$, $\frac{\partial q}{\partial y} = 1$ and f = qz, $\frac{\partial f}{\partial q} = z$, $\frac{\partial f}{\partial z} = q$. We want to calculate $\frac{\partial f}{\partial x}$, $\frac{df}{dy}$, $\frac{\partial f}{\partial z}$.

$$\frac{\partial f}{\partial z} = q = 3$$

$$\frac{\partial f}{\partial q} = z = -4$$

$$\frac{\partial f}{\partial y} = \frac{\partial f}{\partial q} \frac{\partial q}{\partial y} = -4 \cdot 1 = -4$$

$$\frac{\partial f}{\partial x} = \frac{df}{dq} \frac{\partial q}{\partial x} = -4 \cdot 1 = -4$$

9.5 Universal Representation

Theorem (Universal Approximation Theorem by 2-Layer Neural Networks)

For any continuous function $f: \mathbb{R}^d \to \mathbb{R}^c$ and any $\epsilon > 0$, there exists $k \in \mathbb{N}, W \in \mathbb{R}^{k \times d}, \mathbf{b} \in \mathbb{R}^k, U \in \mathbb{R}^{c \times k}$ such that

$$\sup_{\mathbf{x}} \|f(\mathbf{x}) - g(\mathbf{x})\|_2 < \epsilon$$

where $g(\mathbf{x}) = U(\sigma(W\mathbf{x} + \mathbf{b}))$ and σ is the element-wise ReLU operation.

As long as 2-layer MLP is wide enough (large k), it can approximate any continuous function arbitrarily closely.

Why Deep Learning: There exists a function such that a 2-layer MLP needs to be exponentially wide to approximate this function, whereas a 3-layer MLP only needs to be polynomial wide. Deep neural networks are more parameter efficient.

Theorem (Addition is the Only Continuous Multivariate Function)

For any $d \in \mathbb{N}$, there exist constants $\lambda_j > 0$, for j = 1, ..., d, $\mathbf{1}^T \lambda < 1$ and strictly increasing Lipschitz continuous functions $\varphi_k : [0,1] \to [0,1]$ for k = 1, ..., 2d + 1 such that for any continuous function $f : [0,1]^d \to \mathbb{R}$, there exists some continuous univariate function $\sigma : [0,1] \to \mathbb{R}$ so that

$$f(\mathbf{x}) = \sum_{k=1}^{2d+1} \sigma \left(\sum_{j=1}^{d} \lambda_j \varphi_k(x_j) \right)$$

9.5.1 Dropout

For each training minibatch, keep each hidden unit with probability q. We want a different and random network for each training minibatch.

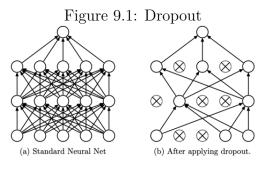
Hidden units are less likely to collude to overfit training data.

Inverted: multiply each \mathbf{h}_k with a scaled binary mask \mathbf{m}/q . We use the full network for testing.

9.5.2 Batch Normalization

Definition: Batch Normalization

Normalize across a minibatch, independently for each feature.



- 1: **Input**: Values of x over a minibatch $\mathcal{B} = \{x_1...\}$ Parameters to be learned: γ, β

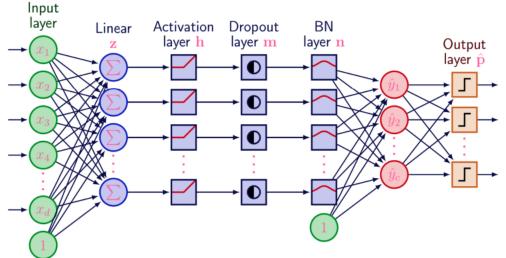
- 2: **Output**: $\{y_i = BN_{\gamma,\beta}(x_i)\}$ 3: $\mu_{\mathcal{B}} \leftarrow \frac{1}{m} \sum_{i=1}^{m} x_i$ 4: $\sigma_{\mathcal{B}}^2 \leftarrow \frac{1}{m} \sum_{i=1}^{m} (x_i \mu_{\mathcal{B}})^2$ 5: $\hat{x}_i \leftarrow \frac{x_i \mu_{\mathcal{B}}}{\sqrt{\sigma_{\mathcal{B}}^2 + \epsilon}}$
- 6: $y_i \leftarrow \gamma \hat{x}_i + \beta \equiv BN_{\gamma,\beta}(x_i)$

⊳ minibatch mean \triangleright minibatch variance \triangleright normalize

Definition: Layer Normalization

Normalize across features, independently for each sample.

Figure 9.2: Complete Multilayer Perceptron Architecture



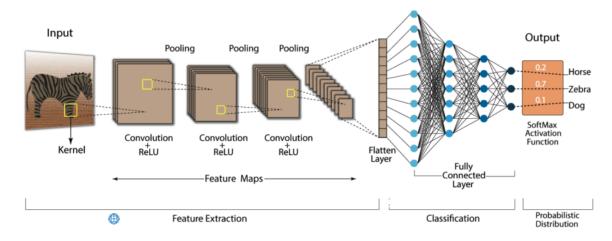
Convolutional Neural Networks

The multilayer perceptron had function $f(\mathbf{x}) = W_2 \sigma(W_1 \mathbf{x} + \mathbf{b}_1) + \mathbf{b}_2$. It had dense weights; each connection represents a weight to be learned, however, it is easy to overfit the training data.

We can consider weight sharing and a sparse weight matrix. We use the same weight for more than one connection and connect less weights.

Layers in a Convolutional Neural Network

Convolution Neural Network (CNN)



Forms of Image Data

- Original image
- Red-green-blue channels: each channel are the RGB images
- Greyscale image

In Pytorch, we represent as a *tensor*. For coloured images, (b, c, m, n) to represent the tensor (batch, channels, rows, columns). For greyscale images, (b, 1, m, n) or (b, m, n).

10.1 Convolution

Definition: Convolution

Let A be the convolutional kernel and X be the input matrix, then the convolution A * X has output size A where each element of A * X is the inner product of A with the |A| submatrix of X.

Definition: One Channel Input

Take the inner product of the convolutional filter with the same size of the input matrix, which will be one element of the output matrix.

The human visual system uses the convolution operation. Traditional image processing algorithms also use the convolution operation.

Definition: Multichannel Input

Take the inner product of each convolutional kernel with the input channel and sum them all for each channel with a bias.

Definition: Filter (Kernel) Size

Width \times height. By default, the number of channels of each filter is the same as that of that of the input.

Definition: Number of Kernels

Weights are not shared between different filters; determine the number of channels of output.

The number of weights to be learned is $c_{in} \times \text{kernel size} \times c_{out}$.

Definition: Stride

Number of pixels to move the filter each time.

Typical stride \leq filter size as to leave no gap. A larger stride makes neighbouring outputs less similar due to less overlap in the input window.

Definition: Padding

Add zeros (or other values) around boundary of input.

10.2 Size Calculation

Input size: $m \times n \times c$

Filter size: $a \times b \times c$

Stride: $s \times t$

Padding: $p \times q$

Pad p pixels on left/right and q pixels on top/bottom (typically p = q).

Move s pixels horizontally and t pixels vertically (typically s = t).

Output size: $\left|1 + \frac{m+2p-a}{s}\right| \times \left|1 + \frac{n+2q-b}{t}\right|$.

With $p = \left\lceil \frac{m(s-1)+a-s}{2} \right\rceil$ and $q = \left\lceil \frac{n(t-1)+b-t}{2} \right\rceil$, the output size = input size.

10.3 Convolutional = FC Layer with Weight Sharing

$$W = \begin{bmatrix} w_{00} & w_{01} \\ w_{10} & w_{11} \end{bmatrix} \in \mathbb{R}^{2 \times 2}, X = \begin{bmatrix} x_{00} & x_{01} & x_{02} \\ x_{10} & x_{11} & x_{12} \\ x_{20} & x_{21} & x_{22} \end{bmatrix} \in \mathbb{R}^{3 \times 3}$$

The circulant matrix is

$$W_{circ} = \begin{bmatrix} w_{00} & w_{01} & 0 & w_{10} & w_{11} & 0 & 0 & 0 & 0 \\ 0 & w_{00} & w_{01} & 0 & w_{10} & w_{11} & 0 & 0 & 0 \\ 0 & 0 & 0 & w_{00} & w_{01} & 0 & w_{10} & w_{11} & 0 \\ 0 & 0 & 0 & 0 & w_{00} & w_{01} & 0 & w_{10} & w_{11} \end{bmatrix} \in \mathbb{R}^{4 \times 9}$$

$$Vector(X) = [x_{00}, x_{01}, x_{02}, x_{10}, x_{11}, x_{12}, x_{20}, x_{21}, x_{22}]^T \in \mathbb{R}^9$$

The circulant matrix W_{circ} is a modification of W where the convolution of W and X will be equal to the matrix multiplication W_{circ} and Vector(X).

$$Vector(W * X) = W_{circ} Vector(X) \in \mathbb{R}^4$$

Suppose we had a fully connected network with size 9 to size 4, then there are $9 \times 4 = 36$ parameters to be learned, but with weight sharing, we only have 4 parameters to learn (the output size).

10.4 Pooling

Definition: Pooling

Down sample input size to reduce computation and memory.

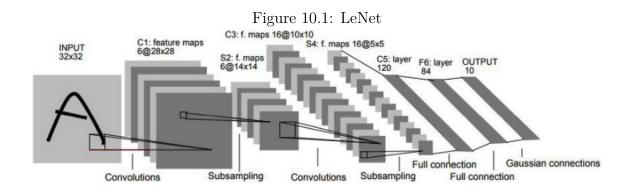
By default, it is performed on each slice separately. Hence output depth equals input depth. There are max-pool, average-pool, (averaged) l_p -norm pool.

Size and stride as in convolution; no parameter; typically no padding.

Definition: Global Pooling

Pooling size = input size.

10.5 CNN Architectures



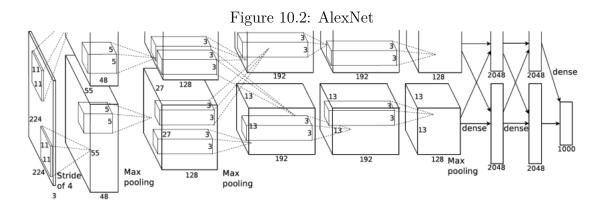


Figure 10.3: LeNet vs. AlexNet √ sigmoid Pool with 2×2 average kernel+2 stride:14×14×6 Pool with 3×3 max. kernel+2 stride: 26×26×96 Pool with 2×2 average kernel+2 stride: 5×5×16 Pool with 3×3 max. kernel+2stride: 12×12×256 √flatten √ sigmoid √ ReLu √ ReLu √ sigmoid ReLu Output: 1 of 10 classes Pool with 3×3 max. kernel+2stride:5×5×256 √ flatten Dense: 4096 fully connected neurons ReLu, dropout p=0.5 Dense: 4096 fully connected neurons ReLu, dropout p=0.5 Dense: 1000 fully connected neurons Output: 1 of 1000 classes

10.5.1 Descriptions

Definition: 1×1 Convolution

Allows easy control of the depth (channels). It can save computation without affecting output size.

GoogLeNet: There are no fully connected (FC) layers. Deeper, but more efficient and better performance.

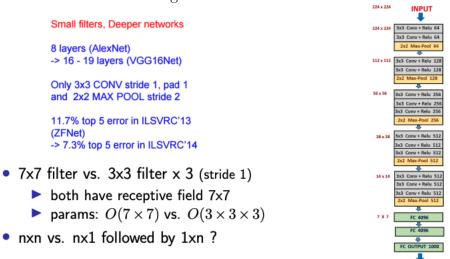
Deeper models can be harder to train due to vanishing/exploding gradient. This can be worse than shallower networks if not properly trained.

Residual Block: We can add a shortcut connection that allows skipping one or more layers. This effectively turns the block into learning residuals: output — input. It allows more direct backpropagation of the gradient through the "shortcut". We can also concatenate or add a linear layer if dimensions mismatch.

ResNet: The full ResNet architecture:

- Stack residual blocks.
- Every residual block has two 3×3 convolutional layers.
- Periodically, double number of filters and down sample spatially using stride 2.
- Additional convolutional layer at the beginning.
- No fully connected layers at the end (only fully connected 1000 to output classes).

Figure 10.4: VGGNet



Filter concatenation

Txt convolutions

3x3 convolutions

3x3 convolutions

3x3 convolutions

3x3 max pooling

1x1 convolutions

1x2 convolutions

1x3 max pooling

Previous layer

(b) Inception module with dimension reductions

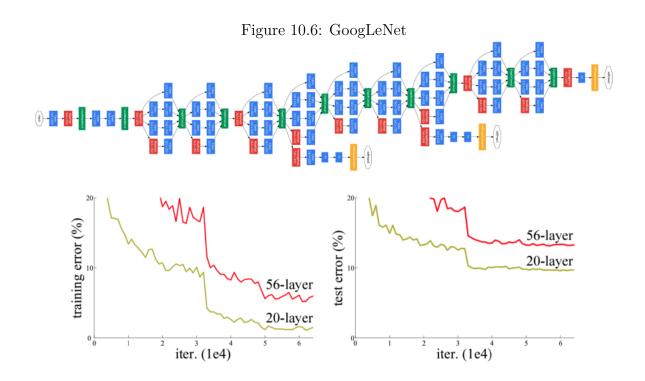


Figure 10.7: Residual Block weight layer $\mathcal{F}(\mathbf{x})$ weight layer \mathbf{x} identity $\mathcal{F}(\mathbf{x}) + \mathbf{x}$ relu

Transformer

Capabilities of transformers: writing code and homework, analyzing texts, taking exams, reasoning, interacting with the real word, etc.

Papers to read:

- Transformer: Attention Is All You Need
- GPT: Improving Language Understanding by Generative Pre-training
- BERT: Pre-training of Deep Bidirectional Transformer for Language Understanding
- GPT-2: Language Models are Unsupervised Multitask Learners
- GPT-3: Language Models are Few-Shot Learners
- GPT-3.5: Training Language Models to Follow Instructions with Human Feedbacks
- GPT-4: Technical Report

Definition: Transformer

Designed for sequence-to-sequence text tasks (e.g. machine translation, question-answer, etc.). The mathematical goal is to output the probability

$$\underset{Y}{\operatorname{arg\,max}} p(\mathbf{y}_1, \dots, \mathbf{y}_m | \mathbf{x}_1, \dots, \mathbf{x}_n)$$

Input and Output

Input Sequence: $X = (\mathbf{x}_1, \dots, \mathbf{x}_n)$, the prompt.

Output Sequence: $Y = (\mathbf{y}_1, \dots, \mathbf{y}_m)$.

Goal: Compute $\arg \max_{Y} p(\mathbf{y}_1, \dots, \mathbf{y}_m | \mathbf{x}_1, \dots, \mathbf{x}_n)$.

encoder

Add & Norm

Feed
Forward

Add & Norm

Multi-Head
Attention

Positional
Encoding

Input
Embedding

Add Milti-Head
Attention

Add & Norm

Masked
Multi-Head
Attention

Positional
Encoding

Figure 11.1: Transformer Architecture

Finding the total output probability might be hard, so we do one token after another (greedy): $\arg\max_{\mathbf{y}_k} p(\mathbf{y}_k|\mathbf{x}_1,\ldots,\mathbf{x}_n,\mathbf{y}_1,\ldots,\mathbf{y}_{k-1})$. This is also known as *auto-regressive*.

Outputs

Example:

- 0. X: Where is University of Waterloo?
- 1. Y: [START]; P(It|X[START]) highest
- 2. Y: [START] It; P(is|X | START] It) highest
- 3. Y: [START] It is; P(in|X [START] It is) highest
- 4. Y: [START] It is in; P(Waterloo|X|START|It is in) highest

Inputs

- 5. Y: [START] It is in Waterloo; P([END]|X | START] It is in Waterloo) highest
- 6. Y: [START] It is in Waterloo [END]

11.1 Input/Output Embedding Layer

Transformer Architecture: In the kth step, the input is x_1, \ldots, x_n and the outputs (shifted right) is y_1, \ldots, y_{k-1} . The output probability is $p(y_k|x_1, \ldots, x_n, y_1, \ldots, y_{k-1})$.

Uses a tokenizer to break up text into individual words. *tiktoken* is a fast tokenizer used by OpenAI's models.

Definition: Token Embedding

A bijection from tokens to vectors.

Convert input tokens to vectors of dimension d = 512. Convert the decode outputted vectors to tokens.

Good properties: words of similar meaning are close in the embedding space.

11.2 Positional Encoding

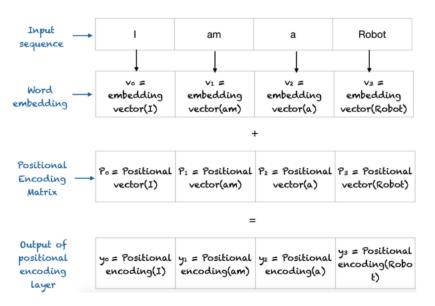
Word order matters, they produce different meanings.

Definition: Positional Encoding

 $W^p \in \mathbb{R}^{n \times d}$.

$$W_{t,2i}^p = \sin(t/10000^{2i/d}), W_{t,2i+1}^p = \cos(t/10000^{2i/d}), i = 0, \dots, \frac{d}{2} - 1$$

There are **NO** parameters to be learned. Simply add W^p to the $n \times d$ token embedding.



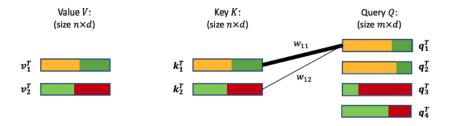
11.3 Attention Layer

Input and Output

Input: Value $V \in \mathbb{R}^{n \times d}$, Key $K \in \mathbb{R}^{n \times d}$, Query $Q \in \mathbb{R}^{m \times d}$.

Output: $m \times d$ matrix.

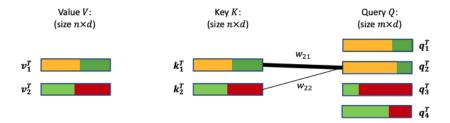
In the case where V=K=Q, they are the output of the positional encoder, and for $V=K\neq Q,\,V,K$ are from the encoder and Q is from the decoder.



 $w_{11} = \langle \mathbf{q}_1, \mathbf{k}_1 \rangle, w_{12} = \langle \mathbf{q}_1, \mathbf{k}_2 \rangle,$ the first output row is

Softmax
$$\left(\frac{w_{11}}{\sqrt{d}}\right)\mathbf{v}_1^T + \operatorname{Softmax}\left(\frac{w_{12}}{\sqrt{d}}\right)\mathbf{v}_2^T$$

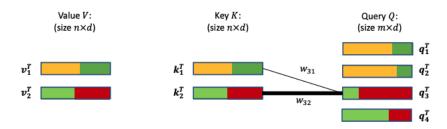
where \mathbf{v}_1^T contributes more to the output row.



 $w_{21} = \langle \mathbf{q}_2, \mathbf{k}_1 \rangle, w_{22} = \langle \mathbf{q}_2, \mathbf{k}_2 \rangle,$ the second output row is

Softmax
$$\left(\frac{w_{21}}{\sqrt{d}}\right)\mathbf{v}_1^T + \operatorname{Softmax}\left(\frac{w_{22}}{\sqrt{d}}\right)\mathbf{v}_2^T$$

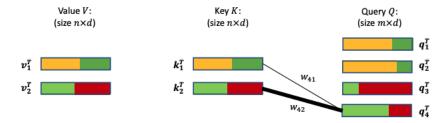
where \mathbf{v}_1^T contributes more to the output row.



 $w_{31} = \langle \mathbf{q}_3, \mathbf{k}_1 \rangle, w_{32} = \langle \mathbf{q}_3, \mathbf{k}_2 \rangle$, the third output row is

Softmax
$$\left(\frac{w_{31}}{\sqrt{d}}\right)\mathbf{v}_1^T + \text{Softmax}\left(\frac{w_{32}}{\sqrt{d}}\right)\mathbf{v}_2^T$$

where \mathbf{v}_2^T contributes more to the output row.



 $w_{41} = \langle \mathbf{q}_4, \mathbf{k}_1 \rangle, w_{42} = \langle \mathbf{q}_4, \mathbf{k}_2 \rangle,$ the fourth output row is

Softmax
$$\left(\frac{w_{41}}{\sqrt{d}}\right)\mathbf{v}_1^T + \text{Softmax}\left(\frac{w_{42}}{\sqrt{d}}\right)\mathbf{v}_2^T$$

where \mathbf{v}_2^T contributes more to the output row.

11.3.1 Matrix Form of Attention

Let $\mathbf{v}_i^T, \mathbf{k}_i^T, \mathbf{q}_i^T$ stand for the row vectors of value, key, and query. Let

$$V = \begin{bmatrix} \mathbf{v}_1^T \\ \vdots \\ \mathbf{v}_n^T \end{bmatrix} \in \mathbb{R}^{n \times d}, K = \begin{bmatrix} \mathbf{k}_1^T \\ \vdots \\ \mathbf{k}_n^T \end{bmatrix} \in \mathbb{R}^{n \times d}, Q = \begin{bmatrix} \mathbf{q}_1^T \\ \vdots \\ \mathbf{q}_m^T \end{bmatrix} \in \mathbb{R}^{m \times d}$$

Then using Softmax (row-wise),

$$\operatorname{Attention}(V, K, Q) = \operatorname{Softmax}\left(\frac{QK^{T}}{\sqrt{d}}\right)V$$

$$= \operatorname{Softmax}\left(\frac{\begin{bmatrix} \mathbf{q}_{1}^{T} \\ \vdots \\ \mathbf{q}_{m}^{T} \end{bmatrix} \begin{bmatrix} \mathbf{k}_{1} & \cdots & \mathbf{k}_{n} \end{bmatrix}}{\sqrt{d}} \right) \begin{bmatrix} \mathbf{v}_{1}^{T} \\ \vdots \\ \mathbf{v}_{n}^{T} \end{bmatrix}$$

$$= \begin{bmatrix} \operatorname{Softmax}\left(\frac{\langle \mathbf{q}_{1}, \mathbf{k}_{1} \rangle}{\sqrt{d}}\right) \mathbf{v}_{1}^{T} + \cdots + \operatorname{Softmax}\left(\frac{\langle \mathbf{q}_{1}, \mathbf{k}_{n} \rangle}{\sqrt{d}}\right) \mathbf{v}_{n}^{T} \\ \vdots \\ \operatorname{Softmax}\left(\frac{\langle \mathbf{q}_{m}, \mathbf{k}_{1} \rangle}{\sqrt{d}}\right) \mathbf{v}_{1}^{T} + \cdots + \operatorname{Softmax}\left(\frac{\langle \mathbf{q}_{m}, \mathbf{k}_{n} \rangle}{\sqrt{d}}\right) \mathbf{v}_{n}^{T} \end{bmatrix} \in \mathbb{R}^{m \times d}$$

$$: \left[\operatorname{Softmax}\left(\frac{\langle \mathbf{q}_{m}, \mathbf{k}_{1} \rangle}{\sqrt{d}}\right) \mathbf{v}_{1}^{T} + \cdots + \operatorname{Softmax}\left(\frac{\langle \mathbf{q}_{m}, \mathbf{k}_{n} \rangle}{\sqrt{d}}\right) \mathbf{v}_{n}^{T} \right] \right]$$

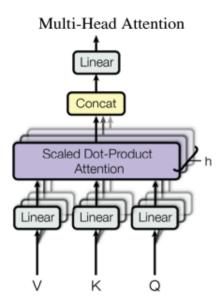
The inner product $\langle \mathbf{q}_i, \mathbf{k}_j \rangle$ measures *similarity*; the more similar, the more contributions. Each output is a convex combination of value rows.

Definition: Self-Attention
$$Q = K = V \label{eq:Q}$$

There is **NO** learnable parameters in this layer.

11.4 Learnable Attention Layer and Multihead Attention

$$\operatorname{Attention}(VW^{v}, KW^{k}, QW^{q}) = \operatorname{Softmax}\left(\frac{QW^{q}(KW^{k})^{T}}{\sqrt{d}}\right)VW^{v}$$



We replace Q, K, V with $QW^q \in \mathbb{R}^{m \times d_k}, KW^k \in \mathbb{R}^{m \times d_k}, VW^v \in \mathbb{R}^{m \times d_v}$ respectively. The learnable linear layers are $\{W^q, W^k, W^v\} \in \mathbb{R}^{512 \times 64}$.

We can add h=8 linear layers W_i^q 's, W_i^k 's, and W_i^v 's in parallel and concatenate their output later. The output dimension is $64 \times 8 = 512$.

11.4.1 Masked Multihead Attention

Do not look at future words, instead mask them. Input the masked sequence into the attention layer.

11.5 Feed-Forward Layer

Definition: Feed-Forward Network

A two-layer MLP with ReLU activation:

$$\mathrm{MLP}(\mathbf{x}) = \max(0, \mathbf{x}^T W_1 + \mathbf{b}_1^T) \cdot W_2 + \mathbf{b}_2^T$$

11.6 Layer Normalization

Batch size (number of sequences) is often small, usually 1 in natural language processing (NLP) tasks. Therefore, batch normalization might not be a good choice. We often use layer normalization instead.

11.7 Overview

There are only 3 tunable hyperparameters. The number of layers N=6, output dimension of all modules is d=512, and the number of attention heads h=8.

The module that connects encoder and decoder is a multihead attention, where the value and key are from encoder, and query from the decoder $(V = K \neq Q)$. In the other two attention modules V = K = Q.

11.8 Transformer Loss

Pretraining task: predict next words. We train by minimizing the log-loss between true next word and predicted next word:

$$\min_{W} \hat{\mathbb{E}}\left[-\left\langle Y, \log \hat{Y} \right\rangle\right]$$

where $Y = [\mathbf{y}_1, \dots, \mathbf{y}_l]$ is the one-hot encoded output sequence and $\hat{Y} = [\hat{\mathbf{y}}_1, \dots, \hat{\mathbf{y}}_l]$ is the predicted probabilities.

Part III

Modern Machine Learning Paradigms

Large Language Models

12.1 Transformer Review

$$\operatorname{Attention}(V, K, Q) = \operatorname{Softmax}\left(\frac{QK^T}{\sqrt{d}}\right)V$$

Each output is a convex combination of value rows. Self-attention means Q = K = V.

The feed-forward network $MLP(\mathbf{x}) = \max(0, \mathbf{x}^T W_1 + \mathbf{b}_1^T) W_2 + \mathbf{b}_2^T$ is a two-layer MLP with ReLU activation with $W_1 \in \mathbb{R}^{d \times d_{ff}}, W_2 \in \mathbb{R}^{d_{ff} \times d}$ where d_{ff} is usually set as 4d.

Attention complexity: $O(n^2d + nd^2)$.

Feed-forward complexity: $O(d^3)$.

where n is the sequence length and d is the dimension of internal representation. Attention also costs $O(nd^2)$ due to linear transformation V = VW where $W \in \mathbb{R}^{d \times d/h}$ where h is the number of heads.

12.2 Labeling Smoothing

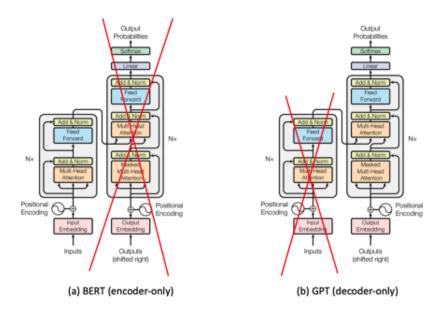
We train large language model (LLM) parameters Θ by minimizing the log-loss between true masked word and predicted masked word:

$$\min_{\Theta} \hat{\mathbb{E}}[-\left\langle Y, \log \hat{Y} \right\rangle]$$

For a training example with ground-truth label y, replace the label Y distribution $p(k|x) = \delta_{k,y}$ with $p'(k|x) = (1 - \epsilon_{ls})\delta_{k,y} + \epsilon_{ls}\frac{1}{C}$ where C is the number of classes.

12.3 BERT vs. GPT

BERT is an encoder, GPT is a decoder. Therefore, BERT is not a generative model.



BERT predicts a randomly-sampled middle word whereas GPT predicts the next word.

BERT is not auto-regressive and GPT is auto-regressive.

12.4 Pretraining-Fine-tuning-Inference Framework

The pretraining stage contains training on large unlabelled datasets such as Wikipedia, Book-Corpus, etc. and is self-supervised training (hours to days). This will learn the pretrained weights.

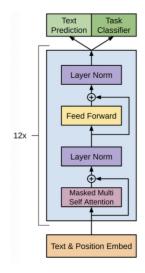
Fine-tuning trains on smaller labelled datasets such as SQuAD, MNLI/CMNLI, Similarity, etc. and is task-specific fine-tuning (minutes to hours). This will fine-tune the weights.

In the pretraining, the model predicts masked words: next words for GPT (harder) or middle words given context for BERT (easier).

It is harder to predict future words, therefore, GPT is better towards artificial general intelligence (AGI).

12.5 Generative Pretraining Transformer (GPT)

The GPT architecture:



12.5.1 Pretraining

Uses unsupervised pretraining through a language model:

$$\min_{\Theta} \hat{\mathbb{E}} \left[\underbrace{-\log \prod_{j=1}^{m} p(\mathbf{x}_{j} | \mathbf{x}_{1}, \dots, \mathbf{x}_{j-1}; \Theta)}_{\text{log-likelihood of predicting next word}} \right]$$

Given the context consisting of previous tokens $\mathbf{x}_1, \dots, \mathbf{x}_{j-1}$ (prompt + previous outputted tokens), predict the current token \mathbf{x}_j .

The pretraining step only uses the Text Prediction output.

12.5.2 Fine-tuning

Uses supervised fine-tuning with task-aware transformations:

$$\min_{\Theta} \underbrace{-\hat{\mathbb{E}}\left[\log\prod_{j=1}^{m}p(\mathbf{y}|X_{1:m};\Theta)\right]}_{\text{task-aware supervised loss}} \cdot -\lambda\hat{\mathbb{E}}\left[\log\prod_{j=1}^{m}p(\mathbf{x}_{j}|X_{1:j-1};\Theta)\right]$$

Fine-tuning uses also uses Task Classifier, the classes are discussed are in the next section.

12.5.3 Task-Dependent Architecture

Definition: Classification

Classify a given text into a class.

Definition: Entailment

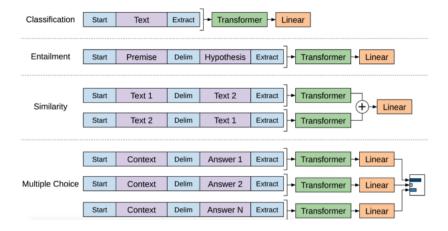
Determine whether a piece of text (hypothesis) contradicts or logically follows from another piece of text (premise).

Definition: Similarity

Predict whether two sentences are semantically equivalent.

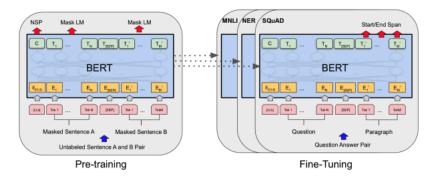
Definition: Multiple Choice

Given a context and N possible answers, choose the correct answer.



Early versions of GPT required task-dependent architectures. More recent version do not require it anymore since models became larger.

12.6 Bidirectional Encoder Representations from Transformers (BERT)



BERT_{BASE} has N = 12, d = 768, h = 12, with total parameters 110M.

BERT_{LARGE} has N=24, d=1024, h=16, with total parameters 340M.

12.6.1 Pretraining Task A

Masked Language Model (Masked LM) randomly selects 15% input tokens to change to [Mask]. It adds softmax to predict the [Mask] tokens. It actually is 12% replaced with [Mask], 1.5% with random.

12.6.2 Pretraining Task B

Next Sentence Prediction (NSP): Given two sentences A and B, 50% of the time B is the actual next sentence that follows A (labeled IsNext), and 50% of the time it is a random sentence from the corpus (labeled NotNext).

The losses for Masked LM and NSP tasks are weighted, summed, and minimized.

 $\mathrm{BERT_{BASE}}$ and GPT have comparable number of parameters, but $\mathrm{BERT_{BASE}}$ has better GLUE score.

12.6.3 RoBERTa

Trained model longer, with bigger batches, over more data. It removed the NSP objective and trained on longer sequences.

BERT uses a cross-encoder: two sentences are passed to the transformer. Given N sentences, there are as many as $\binom{N}{2}$ sentence pairs. 10,000 sentences require about 50M times of inference (65 hours) with BERT.

The sentence-transformer is a twin network that reduces 50M to 10K times of inference and 65 hours to 5 seconds.

12.7 GPT-2

A new dataset of millions of webpages called WebText. It is a 1.5B-parameter transformer with the same training method as GPT. GPT-2 performs on par with BERT on fine-tuning tasks, even though BERT is 10x larger. GPT-2 is good for zero-shot learning.

Definition: Zero-Shot Learning

Model predicts answer given only a natural language description with no gradient updates.

Definition: k-Shot Learning

Model predicts answer given a task description and k examples of the task with no gradient updates.

12.8 GPT-3

A transformer with 175B-parameters. Trained the same as GPT and GPT-2. When the network is large enough, there is in-context learning and chain-of-thought prompting.

In-context learning is few-shot learning.

Chain-of-thought prompting is giving an example of the reasoning process in the prompt. Human labeling chain-of-thought is costly; there is a new research direction prompt: prompy tuning/engineering.

Simply adding "Let's think step by step" in the prompt without a reasoning process can improve performance.

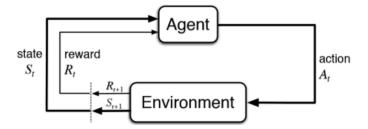
12.9 GPT-3.5

Definition: Policy Function

A function that maps state \rightarrow action.

Uses human feedback to do reinforcement learning. An agent interacts with an environment by taking actions. The environment returns a reward for the action and a new state (representation of world at that moment). Agent uses a policy function to choose an action at a given state.

In the large language model, the state is the prompt, action is the output of the LLM, policy function is the LLM, and the reward is the levels of matching human feedback.



12.10 Reinforcement Learning from Human Feedback

This can be done in 3 steps:

- 1. Update LLM by back-propagation.
- 2. Freeze the LLM and learn the Reward Model (RM).
- 3. Freeze the RM and update the LLM by reinforcement learning.

12.10.1 Step 1

Collect demonstration data and train a supervised policy. Humans provide demonstrations of the desired behaviour on the input prompt distribution. Fine-tune a pretrained GPT-3 model on this data using supervised learning to get the SFT model.

12.10.2 Step 2

Collect comparison data and train a reward model. Collect a dataset of comparisons between model outputs where the labeler indicates which one among K outputs they prefer for a given input. Then train the reward model (RM) to predict human-preferred output.

Pairwise comparison loss to train reward model r_{θ} :

$$loss(\theta) = -\frac{1}{\binom{K}{2}} \mathbb{E}_{(x,y_w,y_l)} [\log(\sigma(r_{\theta}(x,y_w) - r_{\theta}(x,y_l)))]$$

where x is the prompt/input, y_w is the preferred output in (y_w, y_l) , and σ is the sigmoid function. We encourage $r_{\theta}(x, y_w) >> r_{\theta}(x, y_l)$ by logistic loss.

We use the RM model because annotators are hard to give uniformly consistent scores for a given sentence. Since ranking is a discrete reward which suffers from discrete 0-1 reward. The RM model fits the discrete reward and outputs a real-valued reward. This is why we use the RM model instead of a human annotator.

OpenAI has created a human-computer interface for human annotators.

12.10.3 Step 3

Optimize a policy against the reward model using PPO. Use the output of the reward model as a scalar reward. Fine-tune the supervised policy to optimize this reward using a reinforcement learning algorithm.

$$\max_{\phi} \mathbb{E}_{(x,y)} \left[\underbrace{r_{\theta}(x,y)}_{\text{maximize reward}} - \beta \underbrace{\log(\pi_{\phi}^{RL}(y|x)/\pi^{SFT}(y|x))}_{\text{model is close to SFT model}} \right] + \gamma \underbrace{\mathbb{E}[\log(\pi_{\phi}^{RL}(x))]}_{\text{pretraining loss}}$$

ChatGPT and GitHub Copilot are based on GPT-3.5.

12.11 GPT-4

Released on March 14, 2023.

Multi-modality: allows images and text as input and text as output.

12.11.1 Parameter-Efficient Fine-Tuning (PEFT)

Large language models are becoming larger, so the industry has significant advantages over academia.

Definition: Low-Rank Adaptation (LoRA)

Pretrained weights are frozen, only train A and B.

Generative Adversarial Networks

Self-Supervised Learning

Part IV

Trustworthy Machine Learning

Existing Threats – Evasion Attacks

Existing Threats – Physical and Poisoning Attacks

Privacy

Ethics

Robustness