UNSUPERVISED LEARNING

IN5400 — Machine Learning for Image Analysis

Ole-Johan Skrede 03.04.2019

University of Oslo

- · Mandatory 2 is ready soon (some technical difficulties)
- · Exercise for this week is ready before tomorrow

- $\cdot\,$ Introduction and motivation
- \cdot Repetition / background
 - · K nearest neighbours, k-means clustering
 - $\cdot\,$ Principal component analysis
 - · Independent component analysis
- \cdot t-SNE
- · Autoencoders, variational autoencoders

INTRODUCTION AND MOTIVATION

 \cdot Given a training set with pairs of inputs x and corresponding desired outputs y

$$\Omega_{\text{train}} = \{ (x^{(1)}, y^{(1)}), \dots, (x^{(m)}, y^{(m)}) \}$$

 \cdot Create a function f that "approximates" this mapping

 $f(x) \approx y, \quad \forall (x,y) \in \Omega_{\text{train}}$

 \cdot Hope that this generalises well to unseen examples, such that

$$f(x) = \hat{y} \approx y, \quad \forall (x, y) \in \Omega_{\text{test}}$$

where $\Omega_{\rm test}$ is a set of relevant unseen examples.

 \cdot Hope that this is also true for all unseen relevant examples.

- In contrast with supervised learning, we have *no* labeled data points in unsupervised learning.
- Since there is no "ground truth", there is no accuracy evaluation in the supervised sense.
- \cdot Applications
 - · Data clustering
 - · Anomaly detection
 - · Signal generation
 - · Signal compression

- We have *some* labeled data
- $\cdot\,$ Usually a majority of unlabeled data
- $\cdot\,$ Can be thought of as supervised learning extended to utilise unlabeled data
- $\cdot\,$ Will not be covered today

What we will cover today

- · K-means clustering (background)
- Principal component analysis (PCA) (background)
- · t-SNE
- · Autoencoders
- · Variational autoencoders

What we will not cover today

- \cdot Independent component analysis (ICA)
- · Matrix factorization and decomposition
- Expectation-maximization (EM) algorithm
- Generative-adverserial networks (GAN) (next lecture)

CLUSTERING

- \cdot Grouping together data based on some similarity metric
- Data points within the same group (cluster) will be more similar to each other than to data points outside the group
- · Many different versions of clustering

CONNECTIVITY-BASED CLUSTERING

- · Also called hierarchical clustering
- \cdot See figures for example with the L_2 distance metric measured from cluster centroides
- Different level thresholds yields different clusters





Figure 2: Bottom up (agglomerative) hierarchy of clusters

Figure 1: Raw data

GRAPH CLUSTERING — **CLIQUES**

- $\cdot\,$ A clique is a set of nodes
- A node in a clique shares an edge with all other nodes in the clique
- · Can have cliques of different sizes
- $\cdot\,$ Useful in areas such as random fields



Figure 3: Undirected graph



Figure 4: Top: Cliques with 2 members. Bottom: Cliques with 3 members. Nodes with multiple colors belong to more than one clique.

- $\cdot\,$ Clusters are represented by a central vector
- · Example: K-means clustering

K-MEANS CLUSTERING

- · Conseptually simple clustering algorithm
- \cdot We want to partition a set of data $\{x^{(1)}, x^{(2)}, \ldots, x^{(m)}\}$ into k clusters.
- $x^{(i)} \in \mathbb{R}^n, i = 1, \dots, m$
- $\cdot \,$ With some distance norm $|| \cdot ||$ the procedure is
 - 1. Initialize at random k cluster centroids (or means) $\mu_j \in \mathbb{R}^n, j=1,\ldots,k$
 - 2. Repeat until convergence

2.1 Assign every example $x^{(i)}, i=1,\ldots,m$ with the label of the nearest cluster centroid

$$c^{(i)} = \arg\min_{j} ||x^{(i)} - \mu_j||.$$

2.2 Update the position of every centroid $\mu_j, j = 1, \dots, k$ to the centroid of the cluster of points with its label

$$\mu_j = \frac{\sum_{i=1}^m I[c^{(i)} = j]x^{(i)}}{\sum_{i=1}^m I[c^{(i)} = j]},$$

where the Iverson bracket is defined as

$$I[a=b] = \begin{cases} 1, & \text{if } a=b, \\ 0, & \text{if } a\neq b \end{cases}$$

· Minimizes the objective function

$$J(c,\mu) = \sum_{i=1}^{m} ||x^{(i)} - \mu_{c^{(i)}}||$$

- · Not guaranteed to find a *global minimum*
- $\cdot\,$ Common to run the algorighm several times with different initializations, and then pick the run with the smallest value of J
- · The k-means clustering algorithm partitions the feature space into Voronoi cells

K-means clustering — example



Figure 5: A: Initialize centroids. B: Assign points to clusters. C: Move centroids. D, E, F: Assign points to clusters and move centroids. No change after F (convergence).



PRINCIPAL COMPONENT ANALYSIS (PCA)

- · Reducing the dimensionality of a dataset of correlated variables
- \cdot Retaining as much as possible of the variance present in the dataset



Figure 6: Representing 2D data as 1D

- $\cdot \ \operatorname{Let} X \in \mathbf{R}^{n_d}$ be a random vector
- · We are looking for a set of uncorrelated variables Y_k which we will call the principal components of X
- $\cdot\,$ The first component, Y_1 , will account for most of the variance in X
- · The second component, Y_2 , will account for most of the variance in X, conditionied on being uncorrelated with Y_1
- The third component, Y_3 , will account for most of the variance in X, conditioned on being uncorrelated with both Y_1 and Y_2
- $\cdot\,$ We continue untill we have $n_p << n_d$ principal components that account for most of the variance in X

PCA — FIRST PRINCIPAL COMPONENT

 $\cdot \,$ Let $Y_1 \in \mathbb{R}$ be some linear combination of the elements in X

$$Y_1 = \sum_{i=1}^{N_d} a_{1i} X_i = a_1^{\mathsf{T}} X,$$

 \cdot This random variable has variance

$$Var[Y_1] = Var[a_1^{\mathsf{T}}X] = a_1^{\mathsf{T}}\Sigma a_1.$$

 \cdot Here, Σ is the covariance matrix of X with elements

$$\Sigma_{ij} = Cov(X_i, X_j)$$

- \cdot We want to maximize the variance of Y_1
- · In order to achieve finite solutions, we constrain the optimization on

$$a_1^{\mathsf{T}}a_1 = 1$$

· It turns out that, for $k=1,\ldots,n_p$, a_k well be an eigenvector of Σ corresponding to the kth largest eigenvalue λ_k

PCA — ESTIMATING THE COVARIANCE MATRIX

· For a dataset with n_s samples $\{x_{i1}, \ldots, x_{in_s}\}$ for all features $i = 1, \ldots, n_d$, the elements in the covariance matrix can be estimated as

$$\hat{\Sigma}_{ij} = \frac{1}{n_s - 1} \sum_{q=1}^{n_s} (x_{iq} - \hat{\mu}_i)(x_{jq} - \hat{\mu}_j),$$

 \cdot Here $\hat{\mu}_i$ is the sample mean of the *i*th feature

$$\hat{\mu}_i = \frac{1}{n_s} \sum_{q=1}^{n_s} x_{iq}$$

 \cdot Arranginging the feature samples and sample means into vectors of size n_d

$$x_q = [x_{1q}, \dots, x_{n_dq}]^\mathsf{T}$$
$$\hat{\mu} = [\hat{\mu}_1, \dots, \hat{\mu}_{n_d}]^\mathsf{T}$$

 \cdot With this, the estimate of the covariance matrix can be written as

$$\hat{\Sigma} = \frac{1}{n_s - 1} \sum_{q=1}^{n_s} (x_q - \hat{\mu}) (x_q - \hat{\mu})^{\mathsf{T}}.$$

- We use the technique of *Lagrangian multipliers* to incorporate the unit length constraint
- · This means that we are going to maximize the expression

$$J(a_1) = a_1^\mathsf{T} \Sigma a_1 - \lambda (a_1^\mathsf{T} a_1 - 1).$$

 \cdot Computing the gradient of J w.r.t. a_1 , and setting it equal to zero, yields

$$\Sigma a_1 - \lambda a_1 = 0,$$

or

$$(\Sigma - \lambda I)a_1 = 0,$$

where I is the $n_d \times n_d$ identity matrix.

 \cdot From our last expression

$$(\Sigma - \lambda I)a_1 = 0,$$

we see that λ is an eigenvalue of Σ , and a_1 is the corresponding eigenvector.

- \cdot Furthermore, λ is the largest eigenvalue
- This is because maximizing the variance subject to the constraint of unit length coefficients is equivalent to choosing the largest eigenvalue

$$a_1^{\mathsf{T}} \Sigma a_1 = a_1^{\mathsf{T}} \lambda a_1$$
$$= \lambda a_1^{\mathsf{T}} a_1$$
$$= \lambda.$$

 \cdot In general, the *k*th principal component of *X* is

$a_k^{\mathsf{T}} X$

where a_k is the eigenvector of the covariance matrix Σ of X, corresponding to the kth largest eigenvalue λ_k

- · Dimensionality reduction
- · Preprocessing in supervised learning: acts as a regularizer
- \cdot Noise reduction

PROBLEMS WITH IMAGE DATA



(all 3 images have same L2 distance to the one on the left)





MNIST CLUSTERING WITH PCA

Explains about 26% of the variance. Not very suited.





- Precursor to t-SNE (*t-distributed* Stochastic Neighbour Embedding)
- $\cdot\,$ Introduced by Geoffrey Hinton and Sam Roweis in 2003 1
- \cdot A stochastic dimensionality reduction method
- \cdot Transforms high-dimensional (*HD*) data points to low-dimensional (*LD*) data points
- $\cdot\,$ Aims to preserve neighbourhood relationship between data points
- \cdot Similar (close) *HD* points should also be similar (close) in the *LD* representation

¹http://papers.nips.cc/paper/2276-stochastic-neighbor-embedding.pdf

- $\cdot\,$ The high-dimensional points have some dimension h
- \cdot The low-dimensional points have some desired predetermined dimension l << h
- \cdot For each point *i*, we are going to define two distributions:
 - $\cdot P_i(x_j)$: Describes the probability that point j is the "neighbour" of point i, given its location x_i
 - $\cdot \; Q_i(y_j)$: Describes the probability that point j is the "neighbour" of point i, given its location y_i
- \cdot We are then going to define a similarity measure between these distributions
- The low-dimensional representations will be altered such as to minimize this distribution similarity

SNE — HIGH-DIMENSION NEIGHBOUR PROBABILITY

- \cdot Let X be a h-dimensional random variable (RV) modelling a HD point
- $\cdot \,$ Let S be a h-dimensional RV that is modelling a neighbour of X
- · Given that $X = x_i$, we want the probability that S is a neighbour of X to be proportional to the Gaussian of the euclidian distance between the two

$$\Pr(S = s | X = x_i) = \frac{1}{c_i} \frac{1}{\sqrt{2\pi\sigma_i^2}} \exp\left\{-\frac{||x - s||^2}{2\sigma_i^2}\right\}$$

where c_i is a constant

- \cdot We define $\Pr(S = x_i | X = x_i) = 0$
- \cdot We also want it to be a probability, so if we sum over all possible neighbours $z
 eq x_i$

$$\sum_{z \neq x_i} \Pr(S = z | X = x_i) = 1$$

 \cdot We end up with

$$\Pr(S = s | X = x_i) = \frac{\exp\left\{-\frac{||x_i - s||^2}{2\sigma_i^2}\right\}}{\sum_{z \neq x_i} \exp\left\{-\frac{||x_i - s||^2}{2\sigma_i^2}\right\}}$$

SNE — HIGH-DIMENSION NEIGHBOUR PROBABILITY, NOTATION

• The probability mass function that describes the probability that some neighbour S of X is located at s given that X is located at x_i is

$$\begin{split} P_i(s) &:= \Pr(S = s | X = x_i) \\ &:= \frac{\exp\left\{-\frac{||x_i - s||^2}{2\sigma_i^2}\right\}}{\sum_{z \neq x_i} \exp\left\{-\frac{||x_i - z||^2}{2\sigma_i^2}\right\}} \end{split}$$

- \cdot Given a concrete set of points $\{x_1, x_2, \ldots, x_n\}$
- The probability that j is a neighbour of i, given that i is located at x_i is then

$$p_{j|i} := \frac{\exp\left\{-\frac{||x_i - x_j||^2}{2\sigma_i^2}\right\}}{\sum_{k \neq i} \exp\left\{-\frac{||x_i - x_k||^2}{2\sigma_i^2}\right\}}$$

SNE — SCALING PARAMETER

- \cdot The scaling parameter σ_x can be set manually
- \cdot We want a larger σ_x in sparse areas
- \cdot We want a smaller σ_x in dense areas



SNE - PERPLEXITY

- $\cdot \ \sigma$ is often found with binary search such that the perplexity equals k, which is determined manually
- \cdot The perplexity of the distribution P_i is given by

$$Perp(P_i) = 2^{H(P_i)}$$

where the Shannon entropy is given by

$$H(P_i) = -\sum_j p_{j|i} \log_2 p_{j|i}$$

- Perplexity can be interpreted as a measure of how many neighbours we want to influence a point
- Typical values are between 5 and 50
- See e.g. https://distill.pub/2016/misread-tsne/ how to interpret t-SNE results

- $\cdot \,$ Let Y be a l-dimensional RV modelling a LD point
- \cdot Let T be a l-dimensional RV that is modelling a neighbour of Y
- $\cdot \,\, Y$ are the lower-dimensional data points corresponding to X, so l << h
- $\cdot\,$ Similarly to HD , we choose a Gaussian neighbourhood, but with fixed variance $\sigma^2=1/2$

$$\Pr(T = t | Y = y_i) = \frac{\exp\left\{-||y_i - t||^2\right\}}{\sum_{z \neq y_i} \exp\left\{-||y_i - z||^2\right\}}$$

SNE - LOW-DIMENSION NEIGHBOUR PROBABILITY, NOTATION

- \cdot For every HD point x_i , we have a corresponding LD point y_i
- The probability mass function that describes the probability that some neighbour T of Y is located at t given that Y is located at y_i is

$$\begin{aligned} Q_i(s) &:= \Pr(T = t | Y = y_i) \\ &:= \frac{\exp\left\{-||y - t||^2\right\}}{\sum_{z \neq x_i} \exp\left\{-||y_i - z||^2\right\}} \end{aligned}$$

- \cdot Given a concrete set of points $\{y_1, y_2, \dots, y_n\}$
- \cdot The probability that j is a neighbour of i, given that i is located at y_i is then

$$q_{j|i} := \frac{\exp\left\{-||y_i - y_j||^2\right\}}{\sum_{k \neq i} \exp\left\{-||y_i - y_k||^2\right\}}$$

- $\cdot\,$ The goal is to place y_i such that the LD distribution $q_{j|i}$ is similar to the HD distribution $p_{j|i}$
- $\cdot\,$ We need a similarity metric, and a way to optimize it

 \cdot The Kullback-Liebler divergence over a discrete random variable X

$$D_{KL}(p_X||q_X) = \sum_x p_X(x) \log \frac{p_X(x)}{q_X(x)}$$

- Measures the distance between two probability distributions p_X and q_X over the same set of events, modeled with the random variable X.
- \cdot Expectation of logarithmic difference between p and q when expectation is taken w.r.t. p.
- \cdot Measures the amount of information that is lost when using q to approximate p.
- $\cdot\,$ It is non-negative
- · Zero for p = q
- $\cdot\,$ Increasing for "increasing difference" between p and q.
SNE — DISTRIBUTION SIMILARITY MEASURE

- \cdot We want to measure the similarity between P_i and Q_i , for all points i
- This is done by summing the KL-divergence between the original (P_i) and the "induced" (Q_i) distributions over all points

$$C = \sum_{i} D_{KL}(P_i||Q_i)$$
$$= \sum_{i} \sum_{j} p_{j|i} \log \frac{p_{j|i}}{q_{j|i}}$$

- · Large cost of confusing a small distance in the high-dimensional space with a large distance in the low-dimensional space (small $p_{j|i}$ and large $q_{j|i}$)
- · Larger cost of confusing a large distance in the high-dimensional space with a small distance in the low-dimensional space (large $p_{j|i}$ and small $q_{j|i}$)

- \cdot The cost can be minimized with stochastic gradient descent
- · Note that we are minimizing w.r.t. the LD points $\{y_1, \ldots, y_n\}$ corresponding to the known HD points $\{x_1, \ldots, x_n\}$
- \cdot Keeps nearby points in HD nearby in LD
- $\cdot\,$ Also keeps distant points in HD relatively far apart in LD
- $\cdot\,$ Drawback: Can be difficult to optimize
- Drawback: Tendency to crowd *LD* representations at the center of the map ("crowding problem")

- · A variant of the SNE method
- \cdot Introduced by Laurens van der Maaten and Geoffrey Hinton in 2008 ²
- · An improvemet over SNE
 - · Much easier to optimize
 - · Significantly better visualization
- $\cdot\,$ Two major differences between t-SNE and SNE
 - \cdot Symmetric Gaussian point similarity distribution for the HD data points
 - \cdot Student-t point similarity distribution for the LD map points

²https://lvdmaaten.github.io/publications/papers/JMLR_2008.pdf

• Standard SNE use a sum over the KL-divergence between asymmetric conditional probability distributions

$$C = \sum_{i} D_{KL}(P_i || Q_i)$$

- Because of this, different types of errors in the pairwise distances in the map are weighted differently
- · In particular
 - $\cdot\,$ The cost of representing distant data points as close map points is smaller than
 - $\cdot\,$ The cost of representing close data points as distant map points
- \cdot A symmetric cost could ease optimization, and leviate the crowding problem

Symmetric SNE

• In stead, we could use the KL-divergence between symmetric joint probability distributions

$$C = D_{KL}(P||Q)$$
$$= \sum_{i} \sum_{j} p_{ij} \log \frac{p_{ij}}{q_{ij}}$$

 \cdot The joint probability p_{ij} over all points X_i and their neighbours X_j is

$$p_{ij} = \Pr(X_i = x_i, X_j = x_j)$$

- · Again, we define $p_{ii} = 0$, and require that the sum over the entire possibility space (a point and its neighbours for all points) is 1
- $\cdot\,$ With the same Gaussian neighbourhoods as in SNE, we get

$$p_{ij} = \frac{\exp\left\{-\frac{||x_i - x_j||^2}{2\sigma_i^2}\right\}}{\sum_k \sum_{l \neq k} \exp\left\{-\frac{||x_k - x_l||^2}{2\sigma_k^2}\right\}}$$

 \cdot Similarly, for the LD points

$$q_{ij} = \Pr(Y_i = y_i, Y_j = y_j) \\ = \frac{\exp\left\{-\frac{||y_i - y_j||^2}{2\sigma_i^2}\right\}}{\sum_k \sum_{l \neq k} \exp\left\{-\frac{||y_k - y_l||^2}{2\sigma_k^2}\right\}}$$

- \cdot Note that $p_{ij} = p_{ji}$ and $q_{ij} = q_{ji}$
- · Note that this is not what is used in t-SNE, we will come back to that in two slides
- \cdot This is just motivation

Symmetric SNE — high-dimensional space

· We suggested a symmetric, joint probability

$$p_{ij} = \frac{\exp\left\{-\frac{||x_i - x_j||^2}{2\sigma_i^2}\right\}}{\sum_k \sum_{l \neq k} \exp\left\{-\frac{||x_k - x_l||^2}{2\sigma_k^2}\right\}}$$

- · The problem is that for an outlier x_i , $||x_i x_j||$ will be very large (and p_{ij} very small) for all points
- \cdot The placement of the corresponding point y_i will have very little effect on the cost

$$C = \sum_{i} \sum_{j} p_{ij} \log \frac{p_{ij}}{q_{ij}}$$

 \cdot We can fix this by simply using our previous conditional probabilities as

$$p_{ij} = \frac{p_{j|i} + p_{i|j}}{2n}$$

where n is the number of data points

- · With this, we ensure that $\sum_{i} p_{ij} > 1/(2n)$ for all data points x_i
- \cdot Hence, all points x_i are guaranteed to make significant contributions to the cost

- Standard SNE (and other similar methods) suffer from what is known as the crowding problem
- \cdot Too many map points are placed near the center of the map
- \cdot This can be leviated by forcing moderately distant data points to be placed far apart

- To mitigate the crowding problem, we want to give more weight to representing moderately distant data points as close map points
- \cdot The Student-t distribution with one degree of freedom is used

$$q_{ij} = \frac{\left(1 + ||y_i - y_j||^2\right)^{-1}}{\sum_{k \neq l} \left(1 + ||y_k - y_l||^2\right)^{-1}}$$

- · Notice that it is symmetric $q_{ij} = q_{ji}$
- \cdot The Student-t distribution has a much heavier tail than the Gaussian distribution
- $\cdot\,$ Moderate distances in the HD data space are then represented by larger distances in the LD map space

AUTOENCODERS

- An autoencoder is a neural network which purpose is to discover interesting representations of data
- \cdot The idea is to create identity mappings, that is, functions f such that $f(x)\approx x$ for some input x
- It is able to discover interesting representations by enforcing constraints on the network
- \cdot The method requires no labeled data, and is therefore unsupervised

AUTOENCODERS — INTRODUCTION

- $\cdot\,$ An autoencoder consist of an encoder g and an decoder h
- \cdot The encoder maps the input x to some representation z

 $g:x\mapsto z$

 $\cdot\,$ The decoder maps this representation z to some output \hat{x}

 $h:z\mapsto \hat{x}$

 $\cdot\,$ We want to train the encoder and decoder such that

$$f(x) = h(g(x))$$
$$= h(z)$$
$$= \hat{x}$$
$$\approx \hat{x}$$



- $\cdot\,$ Different network constraints leads to different representations z
- · Compression autoencoder
 - $\cdot \,$ If x has d_x dimensions and z has d_z dimensions, and $d_x > d_z$
 - $\cdot\,$ Most common way of constraining the network
- · Denoising autoencoder
 - \cdot Distorting the input x with some random noise
 - $\cdot\,$ Leads to robust representations, resiliant to corrupted input
- Sparse autoencoder
 - $\cdot \, z$ can actually have a greater dimension than x
 - $\cdot\,$ Only allowing a subset of the hidden units to fire at the same time

- \cdot Encoder:
 - · Input -> first hidden layer: fully connected 784 -> 128, relu
 - · 1. hidden -> 2. hidden: fully connected, 128 -> 32, relu
- \cdot Decoder:
 - · 2.hidden -> 3. hidden: fully connected 32 -> 128, relu
 - · 3. hidden -> output: fully connected, 128 -> 784, sigmoid

COMPRESSION AUTOENCODER — MNIST EXAMPLE





- · Same set-up as in a compression autoencoder
- $\cdot\,$ Add noise to the input
- $\cdot\,$ Compare the reconstruction to the input without noise



DENOISING AUTOENCODER – MNIST EXAMPLE

Same setup as for the compression autoencoder. Zero mean Gaussian noise with standard deviation 0.1 is added to the input. The input values are clipped to lay in [0, 1].





- \cdot We want to constrain the number of *active* nodes in the coding layer
- \cdot We can think of a node being active (or *firing*) if is
 - $\cdot\,$ close to 1 for the sigmoid tanh activation functions
- \cdot We can think of a node being inactive
 - $\cdot\,$ close to 0 for the sigmoid activation function
 - $\cdot\,$ close to -1 for the tanh activation function
- \cdot We would like to constrain the nodes to be inactive most of the time

- · Let $a_j^{[c]}(x^{(i)})$ be the activation in node j in the coding layer [c] given an input $x^{(i)}$ to the network
- \cdot Then, activation for this node averaged over all m input examples is

$$\hat{\rho} = \frac{1}{m} \sum_{i=1}^{m} a_j^{[c]}(x^{(i)})$$

 \cdot We would like to limit this average activation by enforcing the constraint

$$\hat{\rho} = \rho$$

for some predetermined sparsity parameter ho

 \cdot Choosing a small ho (e.g. 0.1) forces the activations to be small

 \cdot The way we enforce this constraint is to regularize the loss function

 $L = L_{\rm reconstruction} + \beta L_{\rm sparsity}$

with some regularization strength $\beta \in \mathbb{R}$.

· We are going to use the KL-divergence between the distributions p and q_j summed over the entire latent layer as our sparsity loss

$$L_{\text{sparsity}} = \sum_{j=1}^{n^{[c]}} D_{KL}(p||q_j)$$

where $n^{[c]}$ is the number of nodes in layer [c]

- $\cdot \, \, p$ will be a Bernoulli distribution with mean ρ for a node j
- $\cdot \,\, q_j$ will be a Bernoulli distribution with mean $\hat{
 ho}_j$ for a node j
- The Bernoulli distribution describes the probability of an event with two outcomes (e.g. coin toss)
- · In our case p will represent a node being active with probability ρ , and q_j a node being active with probability $\hat{\rho}_j$

In this case, the KL divergence for a single node j is

$$D_{KL}(p||q_j) = \sum p(x) \log \frac{p(x)}{q_j(x)}$$
$$= \sum p(x) \log p(x) - \sum p(x) \log q_j(x)$$

The support of the distributions is only two outcomes $x \in \{0, 1\}$, and the pmf is

$$p(x) = \begin{cases} (1-\rho), & x = 0 & \text{(the node is inactive)} \\ \rho, & x = 1 & \text{(the node is active)} \end{cases}$$

and conversely for $q_j(x)$. With this, our KL divergence is simply

$$D_{KL}(p||q_j) = \rho \log \rho + (1-\rho) \log(1-\rho) - [\rho \log \hat{\rho} + (1-\rho) \log(1-\hat{\rho}_j)]$$
$$= \rho \log \frac{\rho}{\hat{\rho}_j} + (1-\rho) \log \frac{(1-\rho)}{(1-\hat{\rho}_j)}.$$

 $\cdot\,$ With this, we get our final loss

$$L = L_{\text{reconstruction}} + \beta \sum_{j=1}^{n^{[c]}} \rho \log \frac{\rho}{\hat{\rho}_j} + (1-\rho) \log \frac{(1-\rho)}{(1-\hat{\rho}_j)}$$

- · Remember that $\hat{\rho}_j$ is the *j*th component of $\hat{\rho} = \frac{1}{m} \sum_{i=1}^m a_j^{[c]}(x^{(i)})$
- \cdot This means that we need to average over all examples to compute $\hat{
 ho}$
- $\cdot\,$ This means that we have to encode all said examples
- \cdot In practice, with batch optimization, we average over all examples in a batch

Sparse autoencoder – MNIST example





VARIATIONAL AUTOENCODERS

- · Popular method for signal generation (images, sound, language, etc.)
- · Creating completely new signals
- $\cdot\,$ Or altering existing data
- Especially powerful when you want to alter your data in a specific way, not just randomly

- An autoencoder works great if you want to reconstruct a *replica* of the input
- \cdot Not well suited for generating new signal
- The reason for this is an "unintuitive" latent variable space
- $\cdot\,$ The latent space might be discontinuous
- Random sampling from an "unseen" region of the latent space produces unpredictable results
- No reasonable way to interpolate between categories in the latent space



- A variational autoencoder is designed to have a continuous latent space
- This makes them ideal for random sampling and interpolation
- It achieve this by forcing the encoder g to generate Gaussian representations, $z\sim \mathcal{N}(\mu,\sigma^2)$
- $\cdot\,$ More precisely, for one input, the encoder generates a mean μ and a variance σ^2
- $\cdot \,$ We then sample a zero-mean, unit-variance Gaussian $\tilde{z} \sim \mathcal{N}(0,1)$
- $\cdot \,$ Construct the input z to the decoder from this

$$z = \mu + \tilde{z} \cdot \sigma$$

 $\cdot \,$ With this, z is sampled from $q = \mathcal{N}(\mu, \sigma^2)$



- $\cdot\,$ This is a stochastic sampling
- $\cdot \,$ That is, we can sample different z from the same set of (μ,σ^2)
- The intuition is that the decoder "learns" that for a given input *x*:
 - \cdot the point z is important for reconstruction
 - $\cdot\,$ but also a neighbourhood of z
- In this way, we have smoothed the latent space, at least locally



PROBLEM

- $\cdot \,$ No restriction on μ or σ^2
- $\cdot\,$ Realisticly, clusters of different classes can be placed far apart
- $\cdot\,$ Leaves "empty space" in between with unknown sampling features



- $\cdot\,$ We can guide the solutions by restricting the generative distribution q
- $\cdot\,$ We do this by making it approximate some distribution p
- \cdot In that way, the latent vectors, even for different categories, will be relatively close
- $\cdot\,$ The desired distribution used in variational autoencoders is the standard normal $p=\mathcal{N}(0,1)$
- We use the familiar KL-divergence between the desired and the generated distribution as a regularizer in the loss function
- \cdot With this, the total loss for an example x_i is something like

$$L(x_i) = ||x^{(i)} - f(x^{(i)})|| + D_{KL}(p||q_{\mu_i,\sigma_i})$$

- \cdot That is, the sum of what we call the *reconstruction loss* and the *latent loss*
- \cdot The latent loss for a single input $x^{(i)}$ can be shown to be equal to

$$D_{KL}(p||q_{\mu_i,\sigma_i}) = \frac{1}{2}(\mu_i^2 + \sigma_i^2 - \log \sigma_i^2 - 1)$$

For reference, I will spend some slide deriving the KL Divergence between two Gaussian distributions $p = \mathcal{N}(\mu_p, \sigma_p^2)$ and $q = \mathcal{N}(\mu_q, \sigma_q^2)$. We are going to derive it for the continuous case, where the KL-Divergence can be expressed as

$$D_{KL}(p||q) = \int p(x) \log \frac{p(x)}{q(x)} dx$$
$$= \int p(x) \log p(x) dx - \int p(x) \log q(x) dx$$

We will derive the two terms in the last line seperately

First, for the first term

$$\int p(x) \log p(x) \, \mathrm{d}x = \int p(x) \log \left[(2\pi\sigma_p^2)^{-\frac{1}{2}} \exp\left\{-\frac{(x-\mu_p)^2}{2\sigma_p^2}\right\} \right] \mathrm{d}x$$
$$= -\frac{1}{2} \log(2\pi\sigma_p^2) \int p(x) \, \mathrm{d}x - \frac{1}{2} \int p(x) \frac{(x-\mu_p)^2}{\sigma_p^2} \, \mathrm{d}x \qquad (1)$$
$$= -\frac{1}{2} \log(2\pi\sigma_p^2) - \frac{1}{2\sigma_p^2} \int p(x)(x^2 - 2x\mu_p + \mu_p^2) \, \mathrm{d}x.$$

Similarly, for the second term

$$\int p(x) \log q(x) \, \mathrm{d}x = -\frac{1}{2} \log(2\pi\sigma_q^2) - \frac{1}{2\sigma_q^2} \int p(x) (x^2 - 2x\mu_q + \mu_q^2) \, \mathrm{d}x. \tag{2}$$

KL Divergence between Gaussian distributions

Remember that for a random variable X with pdf f, the expectation is given by

$$E[X] = \int f(x) x \, \mathrm{d}x.$$

Also, we have

$$E[X^{2}] = \int f(x)x^{2} dx$$
$$= Var[X] + E[X]^{2}$$

For the integral in eq. (1), we then get

$$\frac{1}{2\sigma_p^2} \int p(x)(x^2 - 2x\mu_p + \mu_p^2) \, \mathrm{d}x = \frac{1}{2\sigma_p^2} [(\sigma_p^2 + \mu_p^2) - 2\mu_p^2 + \mu_p^2] = \frac{1}{2}.$$
(3)

The integral in eq. (2) is similar,

$$\frac{1}{2\sigma_q^2} \int p(x)(x^2 - 2x\mu_q + \mu_q^2) \, \mathrm{d}x = \frac{1}{2\sigma_q^2} [(\sigma_p^2 + \mu_p^2) - 2\mu_p\mu_q + \mu_q^2] = \frac{\sigma_p^2 + (\mu_p - \mu_q)^2}{2\sigma_q^2}.$$
(4)

KL DIVERGENCE BETWEEN GAUSSIAN DISTRIBUTIONS

Finishing up, using eq. (1) and eq. (2) via eq. (3) and eq. (4), we finally get

$$D_{KL}(p||q) = \int p(x) \log \frac{p(x)}{q(x)} dx$$

= $\int p(x) \log p(x) dx - \int p(x) \log q(x) dx$
= $-\frac{1}{2} \log(2\pi\sigma_p^2) - \frac{1}{2} + \frac{1}{2} \log(2\pi\sigma_q^2) + \frac{\sigma_p^2 + (\mu_p - \mu_q)^2}{2\sigma_q^2}$
= $\frac{1}{2} \left[\log \frac{\sigma_q^2}{\sigma_p^2} + \frac{\sigma_p^2 + (\mu_p - \mu_q)^2}{\sigma_q^2} - 1 \right]$ (5)

When, as in our case $p=\mathcal{N}(\mu,\sigma)$ and $q=\mathcal{N}(0,1)\text{, we get}$

$$D_{KL}(p||q) = \frac{1}{2} \left[\mu^2 + \sigma^2 - \log \sigma^2 - 1 \right].$$

- $\cdot\,$ With a trained variational autoencoder $f=h\circ g$ you can generate new signals
- $\cdot \;$ Sample $z \sim \mathcal{N}_{n^{[c]}}(0,1),$ where $n^{[c]}$ is the number of nodes in the coding layer
- $\cdot \;$ Feed z into the trained decoder h
- $\cdot \ h(z)$ should now be a randomly generated sample from the training distribution


GENERATE NEW SIGNALS - INTERPOLATION

- Say you want to generate a signal *c* that is an interpolation between two signals *a* and *b*
- $\cdot\,$ First, train a variational autoencoder $f=h\circ g$ on the desired distribution
- · Compute mean vectors μ_a and μ_b from encodings g(a) and g(b)
- Compute the average of the two mean vectors

$$\mu_c = \frac{1}{2}(\mu_a + \mu_b)$$

- $\cdot\,$ Then, set the latent variable $z=\mu_c$
- $\cdot c = h(z)$ should then be an interpolation between a and b

Label 1	Label 0	Interpolated
1	D	E
Label 9	Label 4	Interpolated
9	4	9
Label 8	Label 9	Interpolated
8	9	9
Label 9	Label 9	Interpolated
q	9	9

- \cdot Say you want to add a feature of a signal a to the signal b
- $\cdot\,$ You can do this by finding a signal c that is equal to a, except for the specific feature you want
- $\cdot\,$ You can then subtract the latent variable of c from the latent variable of a, and add it to the latent variable of b
- \cdot Then you simply decode the new latent variable
- Example: "Face with glasses = face + (face with glasses face)"
- $\cdot\,$ See examples on the next slides



(a) Interpolation between genders

(b) Add or remove facial features

Figure 11: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 12: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 13: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 14: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 15: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 16: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 17: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 18: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 19: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 20: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 21: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 22: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 23: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 24: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 25: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 26: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 27: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 28: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae



(a) Interpolation between genders

(b) Add or remove facial features

Figure 29: Source: Deep Feature Consistent Variational Autoencoder, https://houxianxu.github.io/assets/project/dfcvae

QUESTIONS?