Advanced Probabilistic Machine Learning SSY316

Linear Models for Classification

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Classification

Goal: Given a training set $\mathcal{D}=\{(\boldsymbol{x}_1,t_1),\ldots,(\boldsymbol{x}_N,t_N)\}$, with $t\in\{C_1,\ldots,C_K\}$, assign (classify) a new example \boldsymbol{x} to one of the classes C_i .

• For convenience: $t \in \{1, \dots, K\}$

Classification divides input space into decision regions \mathcal{R}_i , $i \in [K]$ (points in \mathcal{R}_i assigned to class C_i).

Linear models: Decision boundaries (surfaces) are linear functions of \boldsymbol{x} ((d-1)-dimensional hyperplanes).

Classification: Examples

Data set: annual income, credit card balance, whether or not the person has defaulted

Aim: Predict whether a person will default or not

Data set: Blood tests, X-rays, MRI scans

Aim: Predict whether a person has cancer or not

Classification as a supervised learning problem

Binary classification:

$$egin{aligned} & oldsymbol{x} \in C_1 & \text{if } t = 1 \\ & oldsymbol{x} \in C_2 & \text{if } t = 0 \text{ (or } t = -1) \end{aligned}$$

e.g., $C_1 \equiv \text{cancer}$, $C_2 \equiv \text{no cancer}$

Multi-class classification (K > 2):

Label a point $x \in C_k$ with a $K \times 1$ one-hot vector t with a single one at position k and zeroes elsewhere

$$K=5$$
, class C_2 :

$$t = (0, 1, 0, 0, 0)^{\mathsf{T}}$$

Three distinct modeling approaches

Discriminative deterministic models: Model directly the deterministic mapping between x and t via a parametrized function $\mu(x, w)$ (discriminant function).

More powerful approach: model $p(\mathcal{C}_k|\mathbf{x})$ in an inference stage, then use this distribution to make optimal decisions.

Discriminative probabilistic models: Model $p(C_k|x)$ via a parametrized model.

Generative probabilistic models: Model p(x,t) specifying p(t) ($p(C_k)$) and $p(\boldsymbol{x}|t)$ $(p(\boldsymbol{x}|C_k))$. Then,

$$p(C_k|\boldsymbol{x}) = \frac{p(\boldsymbol{x}|C_k)p(C_k)}{p(\boldsymbol{x})}$$

Why a probabilistic approach?

Compensating for class priors (Cancer):

- Assume prevalence 1/1000
- Naïve classifier that assigns every new case result to "no cancer" class: 99.9% accuracy!
- An unbalanced data set will lead to learning model that likely won't generalize well
- A balanced data set will lead to a more accurate model ... but we need to compensate for modification to training data!

Combining models (Cancer):

- 3 models based on blood tests, X-rays, MRIs
- If models provide posterior probabilities, they can be properly combined,

$$p(C_k|\mathbf{x}_b, \mathbf{x}_x, \mathbf{x}_m) \propto p(\mathbf{x}_b, \mathbf{x}_x, \mathbf{x}_m|C_k)$$

$$= p(\mathbf{x}_b|C_k)p(\mathbf{x}_x|C_k)p(\mathbf{x}_m|C_k)p(C_k)$$

Activation function

Linear regression: $\mu(x, w)$ a linear function of w

• Simplest case: $\mu(\boldsymbol{x}, \boldsymbol{w}) = \boldsymbol{w}^\mathsf{T} \boldsymbol{x} + w_0$

Classification problem: want to predict discrete class labels, or posterior probabilities.

Generalized linear models:

- Generalization of the linear regression model: $\mu(x, w) = a(w^T x + w_0)$
- Decision boundaries: constant $\mu(x) \longrightarrow \text{linear functions of } x!$

Discriminative deterministic models

Discriminative deterministic models: Model directly the deterministic mapping between x and t via a parametrized function $\mu(x, w)$ (discriminant function).

Binary classification:

• $\hat{t}(x)$ obtained by applying a threshold rule on $\mu(x,w)$

Simplest form:

$$\mu(\boldsymbol{x}, \boldsymbol{w}) = a(\boldsymbol{w}^{\mathsf{T}} \boldsymbol{x} + w_0)$$
$$= \boldsymbol{w}^{\mathsf{T}} \boldsymbol{x} + w_0$$

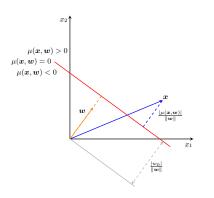
and

$$\hat{t}(\boldsymbol{x}, \boldsymbol{w}) = \operatorname{sign}(\mu(\boldsymbol{x}, \boldsymbol{w}))$$

= $\operatorname{sign}(\boldsymbol{w}^{\mathsf{T}} \boldsymbol{x} + w_0)$
 $\boldsymbol{x} \in \mathcal{C}_1 \text{ if } \mu(\boldsymbol{x}, \boldsymbol{w}) > 0$
 $\boldsymbol{x} \in \mathcal{C}_2 \text{ if } \mu(\boldsymbol{x}, \boldsymbol{w}) < 0$

Geometric interpretation

Decision rule: Defines a hyperplane that separates the domain points classified as belonging to either of the two classes.

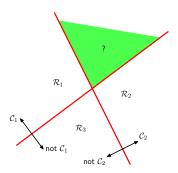


- w determines orientation of decision hyperplane
- w₀ determines location of decision surface
- Classification margin: $|\mu(x, w)| / ||w||$

Combining 2-class discriminant functions:

One-versus the rest classifier: Separating points in C_k from points not in C_k .

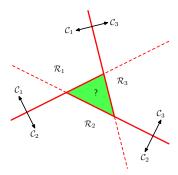
Leads to regions of input space that are ambiguously classified



Combining 2-class discriminant functions:

One-versus-one classifier: $\binom{K}{2}$ binary discriminant functions. Each point classified according to a majority vote.

Leads to regions of input space that are ambiguously classified



A single K-class discriminant with K linear functions

$$\mu_k(\boldsymbol{x}) = w_{k,0} + \boldsymbol{w}_k^\mathsf{T} \boldsymbol{x}$$

and

$$\hat{t}(\boldsymbol{x}) = \arg\max_{k} \mu_k(\boldsymbol{x})$$

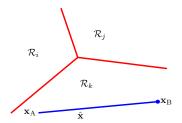
 $\mu_k(\boldsymbol{x})$ can be interpreted as $p(t=k|\boldsymbol{x})$.

Decision boundary between C_k and C_i :

$$\{ \boldsymbol{x} : (\boldsymbol{w}_k - \boldsymbol{w}_j)^\mathsf{T} \boldsymbol{x} + (w_{k,0} - w_{j,0}) = 0 \}$$

A single K-class discriminant with K linear functions

$$\mu_k(\boldsymbol{x}) = w_{k,0} + \boldsymbol{w}_k^\mathsf{T} \boldsymbol{x}$$



- Decision regions singly connected and convex
- Parameters can be optimized via, e.g., least squares $(\boldsymbol{W}^* = (\boldsymbol{X}^\mathsf{T} \boldsymbol{X})^{-1} \boldsymbol{X}^\mathsf{T} \boldsymbol{T})$

Discriminative probabilistic models

Discriminative probabilistic models: Model $p(C_k|x)$ via a parametrized model.

 Potentially more powerful than deterministic models since allow to model sources of uncertainty in the label assignment to input variables

We will consider:

- Models linear in the parameters
- Applying a nonlinear transformation $\phi = (\phi_0, \dots, \phi_M)$ to x, then work with $\phi(x)$ as our features

Decision boundaries linear in the feature space ϕ , nonlinear in original space x.

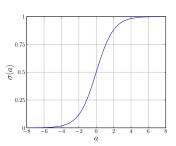
Logistic regression

Binary classification:

$$p(\mathbf{t} = C_1 | \boldsymbol{\phi}, \boldsymbol{x}, \boldsymbol{w}) = p(\mathbf{t} = C_1 | \boldsymbol{\phi}(\boldsymbol{x}), \boldsymbol{w})$$

$$= \frac{\exp(\boldsymbol{w}^\mathsf{T} \boldsymbol{\phi}(\boldsymbol{x}))}{1 + \exp(\boldsymbol{w}^\mathsf{T} \boldsymbol{\phi}(\boldsymbol{x}))}$$

$$= \frac{1}{1 + e^{-(\boldsymbol{w}^\mathsf{T} \boldsymbol{\phi}(\boldsymbol{x}))}} \triangleq \sigma(\boldsymbol{w}^\mathsf{T} \boldsymbol{\phi}(\boldsymbol{x}))$$



For an M-dimensional feature space ϕ , M adjustable parameters.

Logistic regression

$$p(\mathsf{t} = \mathcal{C}_1 | \boldsymbol{\phi}, \boldsymbol{w}) = \frac{1}{1 + e^{-(\boldsymbol{w}^\mathsf{T} \boldsymbol{\phi}(\boldsymbol{x}))}} \triangleq \sigma(\boldsymbol{w}^\mathsf{T} \boldsymbol{\phi}(\boldsymbol{x}))$$

Inference:

Misclassification error minimized by

$$oldsymbol{x} \in \mathcal{C}_1 ext{ if } p(\mathcal{C}_1 | oldsymbol{\phi}, oldsymbol{w})) > 1/2 \ oldsymbol{x} \in \mathcal{C}_2 ext{ if } p(\mathcal{C}_1 | oldsymbol{\phi}, oldsymbol{w})) < 1/2 \$$

Equivalently,

$$oldsymbol{x} \in \mathcal{C}_1 ext{ if } oldsymbol{w}^\mathsf{T} oldsymbol{\phi}(oldsymbol{x}) > 0 \ oldsymbol{x} \in \mathcal{C}_2 ext{ if } oldsymbol{w}^\mathsf{T} oldsymbol{\phi}(oldsymbol{x}) < 0$$

Logistic regression: Learning

ML learning:

$$\boldsymbol{w}^* = \arg\max_{\boldsymbol{w}} \ p(t_{\mathcal{D}}|x_{\mathcal{D}}, \boldsymbol{w})$$

Likelihood function:

$$p(t_{\mathcal{D}}|x_{\mathcal{D}}, \boldsymbol{w}) = \prod_{i=1}^{N} p(t_i|\boldsymbol{x}_i, \boldsymbol{w}) = \prod_{i=1}^{N} p(t_i|\boldsymbol{\phi}(\boldsymbol{x}_i), \boldsymbol{w})$$
$$= \prod_{i=1}^{N} y_i^{t_i} (1 - y_i)^{1 - t_i}$$

with
$$y_i = p(\mathsf{t}_i = \mathcal{C}_1 | \boldsymbol{\phi}(\boldsymbol{x}_i), \boldsymbol{w})$$

Negative log-likelihood:

$$-\ln p(t_{\mathcal{D}}|x_{\mathcal{D}}, \boldsymbol{w}) = -\sum_{i=1}^{N} \left(t_{i} \ln(y_{i}) + (1 - t_{i}) \ln(1 - y_{i}) \right)$$

Logistic regression: Learning

ML learning:

$$w^* = \arg\min_{w} -\ln p(t_{\mathcal{D}}|x_{\mathcal{D}}, w)$$
$$= \arg\min_{w} -\sum_{i=1}^{N} \left(t_i \ln(y_i) + (1 - t_i) \ln(1 - y_i) \right)$$

(cross-entropy loss criterion)

No simple formula to find optimal w^* .

However:

- 1. NLL has unique minimum unless classes perfectly separated by hyperplane
- 2. The negative log-likelihood surface is concave
- 3. The derivatives and the second derivatives can be easily computed

Efficient methods (SGD, Newton-Raphson) exist.

Discriminative probabilistic models

Other discriminative probabilistic models:

- Support vector machine
- Decision tree
- Random forest

We want to obtain

$$p(t|\mathcal{D}, \boldsymbol{x}) = \int \underbrace{p(\boldsymbol{w}|\mathcal{D})}_{\text{posterior dist. of } \boldsymbol{w}} p(t|\boldsymbol{x}, \boldsymbol{w}) d\boldsymbol{w}$$

Bayesian approach: p(t|x, w) associated with each value of w weighted by the posterior belief

$$p(\boldsymbol{w}|\mathcal{D}) = \frac{p(\boldsymbol{w})p(t_{\mathcal{D}}|x_{\mathcal{D}}, \boldsymbol{w})}{p(t_{\mathcal{D}}|x_{\mathcal{D}})}$$

We want to compute the full posterior $p(w|\mathcal{D})$.

Problem: Exact Bayesian inference for logistic regression not possible as exact evaluation of $p(\boldsymbol{w}|\mathcal{D})$ intractable!

$$p(\boldsymbol{w}|\mathcal{D}) \propto p(\boldsymbol{w})p(t_{\mathcal{D}}|x_{\mathcal{D}}, \boldsymbol{w})$$

= $p(\boldsymbol{w}) \prod_{i=1}^{N} p(t_i|\boldsymbol{\phi}(\boldsymbol{x}_i), \boldsymbol{w})$

Bayesian logistic regression: More complex than for linear regression models (cannot integrate exactly over w since posterior not Gaussian) \longrightarrow need to introduce some approximations.

Idea: Use Laplace approximation

The Laplace approximation

Laplace approximation: Aims to find a Gaussian approximation to a posterior distribution.

Case of a single continuous random variable z:

$$p(z) = \frac{1}{Z}f(z)$$

Z: (unknown) normalization constant ($Z = \int f(z) dz$)

Goal: Find a Gaussian approximation q(z) centered on a mode of p(z).

The Laplace approximation

Goal: Find a Gaussian approximation q(z) centered on a mode of p(z).

1. Find a mode of p(z) $(p'(z_0) = 0)$

Observation: The logarithm of a Gaussian is a quadratic function of the variables.

2. Taylor series expansion of $\ln f(z)$ centered on mode z_0 :

$$\ln f(z) \approx \ln f(z_0) + (z - z_0) \frac{d}{dz} \ln f(z) \Big|_{z=z_0} + (z - z_0)^2 \frac{1}{2} \cdot \frac{d^2}{dz^2} \ln f(z) \Big|_{z=z_0}$$

$$= \ln f(z_0) + (z - z_0)^2 \frac{1}{2} \cdot \frac{d^2}{dz^2} \ln f(z) \Big|_{z=z_0}$$

$$= \ln f(z_0) - \frac{1}{2} A(z - z_0)^2$$

with

$$A = -\left. \frac{\mathsf{d}^2}{\mathsf{d}z^2} \ln f(z) \right|_{z=z_0}$$

The Laplace approximation

$$\ln f(z) \approx \ln f(z_0) - \frac{1}{2}A(z - z_0)^2$$

with

$$A = -\left. \frac{\mathsf{d}^2}{\mathsf{d}z^2} \ln f(z) \right|_{z=z_0}$$

We obtain:

$$f(z) \approx f(z_0)e^{-\frac{A}{2}(z-z_0)^2}$$

and normalized distribution

$$q(z) = \left(\frac{A}{2\pi}\right)^{1/2} e^{-\frac{A}{2}(z-z_0)^2}$$
$$= \mathcal{N}(z|z_0, A^{-1})$$

The Laplace approximation for higher dimensions

$$p(z) = \frac{1}{Z}f(z)$$

- 1. Find a mode of p(z) ($\nabla f(z) = 0$)
- 2. Taylor series expansion of $\ln f(z)$ centered on mode z_0 :

$$\ln f(\boldsymbol{z}) pprox \ln f(\boldsymbol{z}_0) - rac{1}{2} (\boldsymbol{z} - \boldsymbol{z}_0)^\mathsf{T} \boldsymbol{A} (\boldsymbol{z} - \boldsymbol{z}_0)$$

with A the $d \times d$ Hessian matrix (evaluated at $z = z_0$)

$$\mathbf{A} = -\nabla\nabla \ln f(z) \Big|_{z=z_0} = \begin{pmatrix} \frac{\partial^2 \ln f(z)}{\partial z_1^2} & \frac{\partial^2 \ln f(z)}{\partial z_1 \partial z_2} & \cdots & \frac{\partial^2 \ln f(z)}{\partial z_1 \partial z_d} \\ \frac{\partial^2 \ln f(z)}{\partial z_2 \partial z_1} & \frac{\partial^2 \ln f(z)}{\partial z_2^2} & \cdots & \frac{\partial^2 \ln f(z)}{\partial z_2 \partial z_d} \\ \vdots & \vdots & \vdots & \vdots \\ \frac{\partial^2 \ln f(z)}{\partial z_d \partial z_1} & \frac{\partial^2 \ln f(z)}{\partial z_d \partial z_2} & \cdots & \frac{\partial^2 \ln f(z)}{\partial z_d^2} \end{pmatrix}$$

The Laplace approximation for higher dimensions

We obtain:

$$f(\mathbf{z}) \approx f(\mathbf{z}_0) e^{-\frac{1}{2}(\mathbf{z} - \mathbf{z}_0)^{\mathsf{T}} \mathbf{A}(\mathbf{z} - \mathbf{z}_0)}$$

and normalized distribution

$$q(z) = \frac{|A|^{1/2}}{(2\pi)^{d/2}} e^{-\frac{1}{2}(z-z_0)^{\mathsf{T}} A(z-z_0)}$$
$$= \mathcal{N}(z|z_0, A^{-1})$$

We want to compute the full posterior $p(w|\mathcal{D})$.

Problem: Exact Bayesian inference for logistic regression not possible as exact evaluation of $p(\boldsymbol{w}|\mathcal{D})$ intractable!

$$p(\boldsymbol{w}|\mathcal{D}) \propto p(\boldsymbol{w})p(t_{\mathcal{D}}|x_{\mathcal{D}}, \boldsymbol{w})$$

= $p(\boldsymbol{w}) \prod_{i=1}^{N} p(t_i|\boldsymbol{x}_i, \boldsymbol{w})$

Idea: Use Laplace approximation

Also need to select a prior for $w \longrightarrow a$ Gaussian prior,

$$p(\boldsymbol{w}) = \mathcal{N}(\boldsymbol{w}|\boldsymbol{m}_0, \boldsymbol{S}_0)$$

Log-likelihood function:

$$\ln p(\boldsymbol{w}|\mathcal{D}) = -\frac{1}{2}(\boldsymbol{w} - \boldsymbol{m}_0)^{\mathsf{T}} \boldsymbol{S}_0^{-1}(\boldsymbol{w} - \boldsymbol{m}_0) + \sum_{i=1}^{N} \left(t_i \ln(y_i) + (1 - t_i) \ln(1 - y_i) \right) + \text{const.}$$

We proceed as follows:

- 1. We maximize $p(w|\mathcal{D})$ (MAP solution w_{MAP}) \longrightarrow defines the mean of the Gaussian
- 2. Covariance given by

$$\boldsymbol{S}_{N}^{-1} = -\nabla\nabla \ln p(\boldsymbol{w}|\mathcal{D}) = \boldsymbol{S}_{0}^{-1} + \sum_{i=1}^{N} y_{i}(1 - y_{i})\boldsymbol{\phi}_{i}\boldsymbol{\phi}_{i}^{\mathsf{T}}$$

Laplace approximation of $p(\boldsymbol{w}|\mathcal{D})$:

$$q(\boldsymbol{w}) = \mathcal{N}(\boldsymbol{w}|\boldsymbol{w}_{\mathsf{MAP}}, \boldsymbol{S}_N)$$

Why Laplace approximation?

 Posteriors more "Gaussian like" as sample-size increases → Good approximation for sufficiently large N

Want to predict $t! \longrightarrow$ Want to learn posterior predictive distribution $p(t|\mathcal{D}, \boldsymbol{x}) = p(t|\mathcal{D}, \boldsymbol{\phi})$

Predictive distribution for class C_1 , $p(t = C_1 | \mathcal{D}, \phi)$:

$$p(C_1|\mathcal{D}, \boldsymbol{\phi}) = \int p(C_1|\boldsymbol{\phi}, \boldsymbol{w}) p(\boldsymbol{w}|\mathcal{D}) d\boldsymbol{w}$$

Given ϕ (or x), we can forecast t via the predictive distribution.

We use

$$p(\mathsf{t} = \mathcal{C}_1 | \boldsymbol{\phi}, \boldsymbol{w}) = \frac{1}{1 + e^{-(\boldsymbol{w}^\mathsf{T} \boldsymbol{\phi}(\boldsymbol{x}))}} \triangleq \sigma(\boldsymbol{w}^\mathsf{T} \boldsymbol{\phi}(\boldsymbol{x}))$$

and

$$q(\boldsymbol{w}) = \mathcal{N}(\boldsymbol{w}|\boldsymbol{w}_{\mathsf{MAP}}, \boldsymbol{S}_N)$$

We obtain:

$$p(C_1|\mathcal{D}, \phi) \approx \int \sigma(\mathbf{w}^\mathsf{T} \phi(\mathbf{x})) q(\mathbf{w}) d\mathbf{w}$$

= $\int \sigma(a) \mathcal{N}(a|\mu_a, \sigma_a^2) da$

with

$$egin{aligned} a &= oldsymbol{w}^\mathsf{T} oldsymbol{\phi} \ \mu_a &= \mathbb{E}[a] = oldsymbol{w}^\mathsf{T}_\mathsf{MAP} oldsymbol{\phi} \ \sigma_a^2 &= \mathsf{Var}[a] = oldsymbol{\phi}^\mathsf{T} oldsymbol{S}_N oldsymbol{\phi} \end{aligned}$$

$$p(C_1|\mathcal{D}, \phi) \approx \int \sigma(a) \mathcal{N}(a|\mu_a, \sigma_a^2) da$$

Cannot be evaluated analytically

Idea: exploit that $\sigma(a)$ is similar to the probit function $\Phi(a)$ to obtain an approximation

$$\Phi(a) \triangleq \int_{-\infty}^{a} \mathcal{N}(x|0,1) dx$$

 $\sigma(a)$ well approximated by $\Phi(\lambda a)$ with $\lambda = \sqrt{\pi/8}$.

$$\int \Phi(\lambda a) \mathcal{N}(a|\mu, \sigma^2) \mathrm{d}a = \Phi\left(\frac{\mu}{(\lambda^{-2} + \sigma^2)^{1/2}}\right)$$

$$\int \Phi(\lambda a) \mathcal{N}(a|\mu, \sigma^2) da = \Phi\left(\frac{\mu}{(\lambda^{-2} + \sigma^2)^{1/2}}\right)$$

Then:

$$p(C_1|\mathcal{D}, \phi) \approx \int \sigma(a) \mathcal{N}(a|\mu_a, \sigma_a^2) da$$
$$\approx \Phi\left(\frac{\mu_a}{(\lambda^{-2} + \sigma_a^2)^{1/2}}\right)$$
$$\approx \sigma\left(\frac{\mu_a}{(1 + \pi \sigma_a^2/8)^{1/2}}\right)$$

Generative probabilistic models

Generative probabilistic models: Model p(x,t) specifying p(t) ($p(C_k)$) and p(x|t) ($p(x|C_k)$). Then,

$$p(C_k|\boldsymbol{x}) = \frac{p(\boldsymbol{x}|C_k)p(C_k)}{p(\boldsymbol{x})}$$

- Make more assumptions about the data by considering distribution of covariates $x \longrightarrow$ may suffer from bias when model incorrectly selected
- Capability to capture properties of x can improve learning if p(x|t) has significant structure

Generative probabilistic models

Binary classification:

$$p(C_1|\mathbf{x}) = \frac{p(\mathbf{x}|C_1)p(C_1)}{p(\mathbf{x}|C_1)p(C_1) + p(\mathbf{x}|C_2)p(C_2)}$$
$$= \frac{1}{1 + e^{-a}} \triangleq \sigma(a)$$

with

$$a = \ln \frac{p(\boldsymbol{x}|\mathcal{C}_1)p(\mathcal{C}_1)}{p(\boldsymbol{x}|\mathcal{C}_2)p(\mathcal{C}_2)}$$

Multi-class classification (K > 2):

$$p(\mathcal{C}_k|\boldsymbol{x}) = \frac{p(\boldsymbol{x}|\mathcal{C}_k)p(\mathcal{C}_k)}{\sum_j p(\boldsymbol{x}|\mathcal{C}_j)p(\mathcal{C}_j)} = \frac{e^{a_k}}{\sum_j e^{a_j}} \qquad \text{(softmax)}$$

with

$$a_k = \ln p(\boldsymbol{x}|C_k)p(C_k)$$

Generative probabilistic models

Model: Different models for $p(x|C_i)$ lead to different generative models. Generative models typically defined by assuming

$$\mathbf{x}|\mathbf{t} = t \sim \text{exponential}(\eta_t)$$

Exponential family:
$$p(\boldsymbol{x}|\boldsymbol{\eta}) = \frac{1}{Z(\boldsymbol{\eta})}h(\boldsymbol{x})\exp(\boldsymbol{\eta}^{\mathsf{T}}\boldsymbol{u}(\boldsymbol{x}))$$

Gaussian discriminant analysis: Class-conditional density modeled by a multivariate Gaussian

- Linear discriminant analysis
- Quadratic discriminant analysis

Linear discriminant analysis

Principle: $x \in \mathbb{R}^d$ is Gaussian distributed for each class C_k , with covariance matrix $\Sigma_k = \Sigma$,

$$\mathbf{x}|\mathbf{t} = C_k \sim \mathcal{N}(\boldsymbol{\mu}_k, \boldsymbol{\Sigma})$$

Equivalently,

$$p(\mathbf{x}|C_k) = \frac{1}{(2\pi)^{d/2} |\mathbf{\Sigma}|^{1/2}} e^{-\frac{1}{2}(\mathbf{x} - \mu_k)^{\mathsf{T}} \mathbf{\Sigma}^{-1} (\mathbf{x} - \mu_k)}$$

Fitting an LDA model to data \equiv estimating Σ and μ_1, \ldots, μ_K .

LDA gives rise to linear decision boundaries.

Linear discriminant analysis

$$\begin{split} \log \frac{p(\mathbf{t} = C_i | \mathbf{x})}{p(\mathbf{t} = C_j | \mathbf{x})} &= \log \frac{p(\mathbf{x} | C_i) p(\mathbf{t} = C_i)}{p(\mathbf{x} | C_j) p(\mathbf{t} = C_j)} \\ &= \log \frac{p(\mathbf{x} | C_i)}{p(\mathbf{x} | C_j)} + \log \frac{p(\mathbf{t} = C_i)}{p(\mathbf{t} = C_j)} \\ &= -\frac{1}{2} \boldsymbol{\mu}_i^\mathsf{T} \boldsymbol{\Sigma}^{-1} \boldsymbol{\mu}_i + \frac{1}{2} \boldsymbol{\mu}_j^\mathsf{T} \boldsymbol{\Sigma}^{-1} \boldsymbol{\mu}_j + \mathbf{x}^\mathsf{T} \boldsymbol{\Sigma}^{-1} (\boldsymbol{\mu}_i - \boldsymbol{\mu}_j) + \log \frac{p(\mathbf{t} = C_i)}{p(\mathbf{t} = C_j)} \\ &= w_0 + \boldsymbol{w}^\mathsf{T} \boldsymbol{x} \end{split}$$

with

$$\begin{aligned} \boldsymbol{w} &= \boldsymbol{\Sigma}^{-1} (\boldsymbol{\mu}_i - \boldsymbol{\mu}_j) \\ \boldsymbol{w}_0 &= -\frac{1}{2} \boldsymbol{\mu}_i^\mathsf{T} \boldsymbol{\Sigma}^{-1} \boldsymbol{\mu}_i + \frac{1}{2} \boldsymbol{\mu}_j^\mathsf{T} \boldsymbol{\Sigma}^{-1} \boldsymbol{\mu}_j + \log \frac{p(C_i)}{p(C_j)} \end{aligned}$$

Decision boundaries: points for which
$$p(\mathbf{t} = C_i | \mathbf{x}) = p(\mathbf{t} = C_j | \mathbf{x}) \longrightarrow \{\mathbf{x} : w_0 + \mathbf{w}^\mathsf{T} \mathbf{x} = 0\}$$

Quadratic discriminant analysis

Principle: $x \in \mathbb{R}^D$ is Gaussian distributed for each class C_k , with covariance matrix \sum_{k}

$$\mathbf{x}|\mathbf{t} = C_k \sim \mathcal{N}(\boldsymbol{\mu}_k, \boldsymbol{\Sigma}_k)$$

Equivalently,

$$p(\boldsymbol{x}|C_k) = \frac{1}{(2\pi)^{D/2} |\boldsymbol{\Sigma}_k|^{1/2}} e^{-\frac{1}{2}(\boldsymbol{x} - \boldsymbol{\mu}_k)^\mathsf{T} \boldsymbol{\Sigma}_i^{-1} (\boldsymbol{x} - \boldsymbol{\mu}_k)}$$

Discriminant functions:

$$\log \frac{p(\mathsf{t} = C_i | \boldsymbol{x})}{p(\mathsf{t} = C_j | \boldsymbol{x})} = -\frac{1}{2} \log |\boldsymbol{\Sigma}_i| - \frac{1}{2} (\boldsymbol{x} - \boldsymbol{\mu}_i)^\mathsf{T} \boldsymbol{\Sigma}_i^{-1} (\boldsymbol{x} - \boldsymbol{\mu}_i) + \log \frac{p(C_i)}{p(C_j)}$$

The decision boundaries are given by a quadratic equation.

Need to fit model: Determine values of w and $p(C_k)$ via maximum likelihood criterion.

Binary case:

- N training examples $\mathcal{D} = \{(x_i, t_i)\}$
- t=1 for class C_1 and t=0 for class C_2
- $p(C_1) = \pi$, $p(C_2) = 1 \pi$

We have:

$$p(\boldsymbol{x}_i, \mathcal{C}_1) = p(\mathcal{C}_1)p(\boldsymbol{x}_i|\mathcal{C}_1) = \pi \mathcal{N}(\boldsymbol{x}_i|\boldsymbol{\mu}_1, \Sigma)$$

$$p(\boldsymbol{x}_i, \mathcal{C}_2) = p(\mathcal{C}_2)p(\boldsymbol{x}_i|\mathcal{C}_2) = (1 - \pi)\mathcal{N}(\boldsymbol{x}_i|\boldsymbol{\mu}_2, \Sigma)$$

Likelihood function:

$$p(\mathcal{D}|\pi, \boldsymbol{\mu}_1, \boldsymbol{\mu}_2, \boldsymbol{\Sigma}) = \prod_{i=1}^{N} \left(\pi \mathcal{N}\left((\boldsymbol{x}_i|\mu_1, \boldsymbol{\Sigma})\right)^{t_i} \left((1-\pi) \mathcal{N}\left((\boldsymbol{x}_i|\mu_2, \boldsymbol{\Sigma})\right)^{1-t_i}\right)\right)^{t_i}$$

Likelihood function

$$p(\mathcal{D}|\pi, \boldsymbol{\mu}_1, \boldsymbol{\mu}_2, \boldsymbol{\Sigma}) = \prod_{i=1}^{N} \left(\pi \mathcal{N}\left((\boldsymbol{x}_i|\boldsymbol{\mu}_1, \boldsymbol{\Sigma})\right)^{t_i} \left((1-\pi) \mathcal{N}\left((\boldsymbol{x}_i|\boldsymbol{\mu}_2, \boldsymbol{\Sigma})\right)^{1-t_i}\right)\right)^{t_i}$$

Log-likelihood function:

$$\begin{split} & \ln p(\mathcal{D}|\pi, \boldsymbol{\mu}_1, \boldsymbol{\mu}_2, \boldsymbol{\Sigma}) \\ & = \sum_{i=1}^{N} \left(t_i \ln \left(\pi \mathcal{N} \left((\boldsymbol{x}_i | \boldsymbol{\mu}_1, \boldsymbol{\Sigma}) \right) + (1 - t_i) \ln \left((1 - \pi) \mathcal{N} \left((\boldsymbol{x}_i | \boldsymbol{\mu}_2, \boldsymbol{\Sigma}) \right) \right) \right) \\ & = \sum_{i=1}^{N} \left(t_i \ln \left(\pi \right) + (1 - t_i) \ln (1 - \pi) \right) + \sum_{i=1}^{N} \left(t_i \ln \left(\mathcal{N} \left(\boldsymbol{x}_i | \boldsymbol{\mu}_1, \boldsymbol{\Sigma} \right) \right) \right) \\ & + \sum_{i=1}^{N} \left((1 - t_i) \ln \left(\mathcal{N} \left(\boldsymbol{x}_i | \boldsymbol{\mu}_2, \boldsymbol{\Sigma} \right) \right) \right) \end{split}$$

Can estimate π , μ_i , and Σ separately!

Maximization with respect to π :

$$\sum_{i=1}^{N} \left(t_i \ln (\pi) + (1 - t_i) \ln (1 - \pi) \right)$$

Differentiating and equating to zero:

$$\frac{\mathsf{d}}{\mathsf{d}\pi} \sum_{i=1}^{N} \left(t_i \ln (\pi) + (1 - t_i) \ln (1 - \pi) \right) = \sum_{i=1}^{N} t_i \frac{1}{\pi} - (1 - t_i) \frac{1}{1 - \pi}$$

$$\implies \pi = \frac{1}{N} \sum_{i=1}^{N} t_i = \frac{N_1}{N}$$

 N_1 : number of data points in C_1

Maximization with respect to μ_1 :

$$\sum_{i=1}^{N} \left(t_i \ln \left(\mathcal{N} \left(\boldsymbol{x}_i | \mu_1, \boldsymbol{\Sigma} \right) \right) \right) = -\frac{1}{2} \sum_{i=1}^{N} t_i (\boldsymbol{x}_i - \boldsymbol{\mu}_1)^\mathsf{T} \boldsymbol{\Sigma}^{-1} (\boldsymbol{x}_i - \boldsymbol{\mu}_1) + \mathsf{const.}$$

Differentiating and equating to zero:

$$\boldsymbol{\mu}_1 = \frac{1}{N_1} \sum_{i=1}^{N} t_i \boldsymbol{x}_i$$

Similarly,

$$\mu_2 = \frac{1}{N_2} \sum_{i=1}^{N} (1 - t_i) x_i$$

Maximization with respect to Σ :

$$oldsymbol{\Sigma} = rac{N_1}{N} oldsymbol{\Sigma}_1 + rac{N_2}{N} oldsymbol{\Sigma}_2$$

with

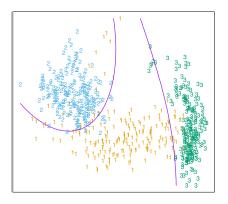
$$oldsymbol{\Sigma}_1 = rac{1}{N_1} \sum_{\substack{i=1:\t_i = C_1}}^{N} (oldsymbol{x}_i - oldsymbol{\mu}_1) (oldsymbol{x}_i - oldsymbol{\mu}_1)^\mathsf{T}$$

$$oldsymbol{\Sigma}_2 = rac{1}{N_2} \sum_{\substack{i=1:\t_i=C_2}}^{N} (oldsymbol{x}_i - oldsymbol{\mu}_2) (oldsymbol{x}_i - oldsymbol{\mu}_2)^{\mathsf{T}}$$

Multi-class classification (K > 2):

$$egin{aligned} \pi_k &= rac{N_k}{N} \ oldsymbol{\mu}_k &= rac{1}{N_k} \sum_{\substack{i=1:\tilde{t}_i = C_k}}^{N} t_i oldsymbol{x}_i \ oldsymbol{\Sigma}_k &= rac{1}{N_k} \sum_{\substack{i=1:\tilde{t}_i = C_k}}^{N} (oldsymbol{x}_i - oldsymbol{\mu}_k) (oldsymbol{x}_i - oldsymbol{\mu}_k)^{\mathsf{T}} \end{aligned}$$

Quadratic discriminant analysis



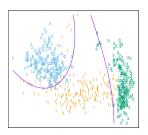
ullet Data $\{oldsymbol{x}_i=(x_{i,1},x_{i,2})\}$ generated from three different Gaussian mixtures

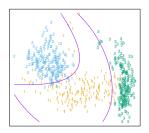
Gaussian discriminant analysis

Two alternative ways to realize quadratic decision boundaries:

- 1. Apply QDA in the original d-dimensional space $x = (x_1, \dots, x_d)$
- 2. Apply LDA in an augmented space of dimension d^2+1 : $x'=(x_1,\ldots,x_d,x_1^2,\ldots,x_d^2,x_1,x_2,\ldots,x_{d-1}x_d)$

Linear functions in the augmented space correspond to quadratic functions in the original *d*-dimensional space.





- Left: QDA applied to original 2-dimensional space $x = (x_1, x_2)$
- Right: LDA applied to 5-dimensional space $x = (x_1, x_2, x_1x_2, x_1^2, x_2^2)$

Gaussian discriminant analysis

Assumption in LDA and QDA: $p(x|C_i)$ are Gaussian distributed ... and yet perform well on a large variety of classification tasks!

Observations:

- From $p(x|\mathcal{C}_k)$ and $p(\mathcal{C}_k)$, we use Bayes' theorem to find $p(\mathcal{C}_k|x)$
- Generative probabilistic model: We can generate synthetic data by drawing values of x from p(x)
- ML learning: fast and simple...but it may overfit in high dimensions

Discriminative or generative models?

Easy to fit?

- Generative: Easy (counting and averaging)
- Discriminative (logistic regression): Solving a complex optimization problem

Adding more classes?

- Generative: No need to retrain (parameters of each class conditional density estimated independently)
- Discriminative: Need to retrain (all parameters interact)

Feature processing?

- Generative: Hard
- Discriminative: Easy, we can simply replace x by $\phi(x)$

Reading

"Pattern recognition and machine learning,"

Chapter 4 (Intro, 4.1.1-4.1.3, 4.1.7, 4.2 (until 4.2.3), 4.3 (until 4.3.4), 4.4. (not 4.4.1), 4.5)