# MACHINE LEARNING

## **Dimensionality reduction**

Corso di Laurea Magistrale in Informatica

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a.a. 2023-2024



In general, many features: high-dimensional spaces.

- sparseness of data
- increase in the number of coefficients, for example for dimension *D* and order 3 of the polynomial,

$$y(\mathbf{x}, \mathbf{w}) = w_0 + \sum_{i=1}^{D} w_i x_i + \sum_{i=1}^{D} \sum_{j=1}^{D} w_{ij} x_i x_j + \sum_{i=1}^{D} \sum_{j=1}^{D} \sum_{k=1}^{D} w_{ijk} x_i x_j x_k$$

number of coefficients is  $O(D^{M})$ 

High dimensions lead to difficulties in machine learning algorithms (lower reliability or need of large number of coefficients) this is denoted as curse of dimensionality

#### **DIMENSIONALITY REDUCTION**

- for any given classifier, the training set size required to obtain a certain accuracy grows exponentially wrt the number of features
- it is important to bound the number of features, identifying the less discriminant ones

#### **DIMENSIONALITY REDUCTION**

- Feature selection: identify a subset of features which are still discriminant, or, in general, still represent most dataset variance
- Feature extraction: identify a projection of the dataset onto a lower-dimensional space, in such a way to still represent most dataset variance
  - Linear projection: principal component analysis, probabilistic PCA, factor analysis
  - Non linear projection: manifold learning, autoencoders

## SEARCHING HYPERPLANES FOR THE DATASET

• verifying whether training set elements lie on a hyperplane (a space of lower dimensionality), apart from a limited variability (which could be seen as noise)



- principal component analysis looks for a d'-dimensional subspace (d' < d) such that the projection of elements onto such suspace is a "faithful" representation of the original dataset
- as "faithful" representation we mean that distances between elements and their projections are small, even minimal

• Objective: represent all *d*-dimensional vectors  $\mathbf{x}_1, \ldots, \mathbf{x}_n$  by means of a unique vector  $\mathbf{x}_0$ , in the most faithful way, that is so that

$$J(\mathbf{x}_0) = \sum_{i=1}^n ||\mathbf{x}_0 - \mathbf{x}_i||^2$$

is minimum

• it is easy to show that

$$\mathbf{x}_0 = \mathbf{m} = \frac{1}{n} \sum_{i=1}^n \mathbf{x}_i$$

• In fact,

 $J(\mathbf{x}_{0}) = \sum_{i=1}^{n} ||(\mathbf{x}_{0} - \mathbf{m}) - (\mathbf{x}_{i} - \mathbf{m})||^{2}$ =  $\sum_{i=1}^{n} ||\mathbf{x}_{0} - \mathbf{m}||^{2} - 2\sum_{i=1}^{n} (\mathbf{x}_{0} - \mathbf{m})^{T} (\mathbf{x}_{i} - \mathbf{m}) + \sum_{i=1}^{n} ||\mathbf{x}_{i} - \mathbf{m}||^{2}$ =  $\sum_{i=1}^{n} ||\mathbf{x}_{0} - \mathbf{m}||^{2} - 2(\mathbf{x}_{0} - \mathbf{m})^{T} \sum_{i=1}^{n} (\mathbf{x}_{i} - \mathbf{m}) + \sum_{i=1}^{n} ||\mathbf{x}_{i} - \mathbf{m}||^{2}$ =  $\sum_{i=1}^{n} ||\mathbf{x}_{0} - \mathbf{m}||^{2} + \sum_{i=1}^{n} ||\mathbf{x}_{i} - \mathbf{m}||^{2}$ 

since

$$\sum_{i=1}^{n} (\mathbf{x}_i - \mathbf{m}) = \sum_{i=1}^{n} \mathbf{x}_i - \mathbf{n} \cdot \mathbf{m} = \mathbf{n} \cdot \mathbf{m} - \mathbf{n} \cdot \mathbf{m} = 0$$

• the second term is independent from  $\mathbf{x}_0$ , while the first one is equal to zero for  $\mathbf{x}_0 = \mathbf{m}$ 



- a single vector is too concise a representation of the dataset: anything related to data variability gets lost
- a more interesting case is the one when vectors are projected onto a line passing through m



• let  $\mathbf{u}_1$  be unit vector ( $||\mathbf{u}_1|| = 1$ ) in the line direction: the line equation is then

#### $\mathbf{x} = \alpha \mathbf{u}_1 + \mathbf{m}$

where  $\alpha$  is the distance of **x** from **m** along the line

• let  $\tilde{\mathbf{x}}_i = \alpha_i \mathbf{u}_1 + \mathbf{m}$  be the projection of  $\mathbf{x}_i$  (i = 1, ..., n) onto the line: given  $\mathbf{x}_1, ..., \mathbf{x}_n$ , we wish to find the set of projections minimizing the quadratic error

The quadratic error is defined as

$$J(\alpha_1, \dots, \alpha_n, \mathbf{u}_1) = \sum_{i=1}^n ||\tilde{\mathbf{x}}_i - \mathbf{x}_i||^2$$
  
=  $\sum_{i=1}^n ||(\mathbf{m} + \alpha_i \mathbf{u}_1) - \mathbf{x}_i||^2$   
=  $\sum_{i=1}^n ||\alpha_i \mathbf{u}_1 - (\mathbf{x}_i - \mathbf{m})||^2$   
=  $\sum_{i=1}^n +\alpha_i^2 ||\mathbf{u}_1||^2 + \sum_{i=1}^n ||\mathbf{x}_i - \mathbf{m}||^2 - 2\sum_{i=1}^n \alpha_i \mathbf{u}_1^{\mathsf{T}}(\mathbf{x}_i - \mathbf{m})$   
=  $\sum_{i=1}^n \alpha_i^2 + \sum_{i=1}^n ||\mathbf{x}_i - \mathbf{m}||^2 - 2\sum_{i=1}^n \alpha_i \mathbf{u}_1^{\mathsf{T}}(\mathbf{x}_i - \mathbf{m})$ 

Its derivative wrt  $\alpha_k$  is

$$\frac{\partial}{\partial \alpha_k} J(\alpha_1, \ldots, \alpha_n, \mathbf{u}_1) = 2\alpha_k - 2\mathbf{u}_1^{\mathsf{T}}(\mathbf{x}_k - \mathbf{m})$$

which is zero when  $\alpha_k = \mathbf{u}_1^{\mathsf{T}}(\mathbf{x}_k - \mathbf{m})$  (the orthogonal projection of  $\mathbf{x}_k$  onto the line).

The second derivative turns out to be positive

$$\frac{\partial}{\partial \alpha_k^2} J(\alpha_1, \ldots, \alpha_n, \mathbf{u}_1) = 2$$

showing that what we have found is indeed a minimum.

To derive the best direction  $\mathbf{u}_1$  of the line, we consider the covariance matrix of the dataset

$$\mathbf{S} = \frac{1}{n} \sum_{i=1}^{n} (\mathbf{x}_i - \mathbf{m}) (\mathbf{x}_i - \mathbf{m})^{\mathsf{T}}$$

By plugging the values computed for  $\alpha_i$  into the definition of  $J(\alpha_1, \ldots, \alpha_n, \mathbf{u}_1)$ , we get

$$\begin{aligned} (\mathbf{u}_1) &= \sum_{i=1}^n \alpha_i^2 + \sum_{i=1}^n ||\mathbf{x}_i - \mathbf{m}||^2 - 2\sum_{i=1}^n \alpha_i^2 \\ &= -\sum_{i=1}^n [\mathbf{u}_1^T (\mathbf{x}_i - \mathbf{m})]^2 + \sum_{i=1}^n ||\mathbf{x}_i - \mathbf{m}||^2 \\ &= -\sum_{i=1}^n \mathbf{u}_1^T (\mathbf{x}_i - \mathbf{m}) (\mathbf{x}_i - \mathbf{m})^T \mathbf{u}_1 + \sum_{i=1}^n ||\mathbf{x}_i - \mathbf{m}||^2 \\ &= -n\mathbf{u}_1^T \mathbf{S} \mathbf{u}_1 + \sum_{i=1}^n ||\mathbf{x}_i - \mathbf{m}||^2 \end{aligned}$$

- $\mathbf{u}_1^T(\mathbf{x}_i \mathbf{m})$  is the projection of  $\mathbf{x}_i$  onto the line
- the product

$$\mathbf{u}_1^{\mathsf{T}}(\mathbf{x}_i - \mathbf{m})(\mathbf{x}_i - \mathbf{m})^{\mathsf{T}}\mathbf{u}_1$$

is then the variance of the projection of  $\mathbf{x}_i$  wrt the mean  $\mathbf{m}$ 

• the sum

$$\sum_{i=1}^{n} \mathbf{u}_{1}^{\mathsf{T}} (\mathbf{x}_{i} - \mathbf{m}) (\mathbf{x}_{i} - \mathbf{m})^{\mathsf{T}} \mathbf{u}_{1} = n \mathbf{u}_{1}^{\mathsf{T}} \mathbf{S} \mathbf{u}_{1}$$

is the overall variance of the projections of vectors  $\mathbf{x}_i$  wrt the mean  $\mathbf{m}$ 

Minimizing  $J(\mathbf{u}_1)$  is equivalent to maximizing  $\mathbf{u}_1^T \mathbf{S} \mathbf{u}_1$ . That is,  $J(\mathbf{u}_1)$  is minimum if  $\mathbf{u}_1$  is the direction which keeps the maximum amount of variance in the dataset

Hence, we wish to maximize  $\mathbf{u}_1^T \mathbf{S} \mathbf{u}_1$  (wrt  $\mathbf{u}_1$ ), with the constraint  $||\mathbf{u}_1|| = 1$ .

By applying Lagrange multipliers this results equivalent to maximizing

 $\boldsymbol{u} = \mathbf{u}_1^\mathsf{T} \mathbf{S} \mathbf{u}_1 - \lambda_1 (\mathbf{u}_1^\mathsf{T} \mathbf{u}_1 - 1)$ 

This can be done by setting the first derivative wrt  $\mathbf{u}_1$ :

$$\frac{\partial \boldsymbol{u}}{\partial \mathbf{u}_1} = 2\mathbf{S}\mathbf{u}_1 - 2\lambda_1\mathbf{u}_1$$

to 0, obtaining

$$\mathbf{S}\mathbf{u}_1 = \lambda_1 \mathbf{u}_1$$



Note that:

- u is maximized if  $\mathbf{u}_1$  is an eigenvector of  $\mathbf{S}$
- the overall variance of the projections is then equal to the corresponding eigenvalue

 $\mathbf{u}_1^{\mathsf{T}} \mathbf{S} \mathbf{u}_1 = \mathbf{u}_1^{\mathsf{T}} \lambda_1 \mathbf{u}_1 = \lambda_1 \mathbf{u}_1^{\mathsf{T}} \mathbf{u}_1 = \lambda_1$ 

• the variance of the projections is then maximized (and the error minimized) if  $\mathbf{u}_1$  is the eigenvector of  $\mathbf{S}$  corresponding to the maximum eigenvalue  $\lambda_1$ 

- The quadratic error is minimized by projecting vectors onto a hyperplane defined by the directions associated to the d' eigenvectors corresponding to the d' largest eigenvalues of S
- If we assume data are modeled by a *d*-dimensional gaussian distribution with mean μ and covariance matrix Σ, PCA returns a *d'*-dimensional subspace corresponding to the hyperplane defined by the eigenvectors associated to the *d'* largest eigenvalues of Σ
- The projections of vectors onto that hyperplane are distributed as a *d'*-dimensional distribution which keeps the maximum possible amount of data variability

# AN EXAMPLE OF PCA

• Digit recognition ( $D = 28 \times 28 = 784$ )



# CHOOSING d'

Eigenvalue size distribution is usually characterized by a fast initial decrease followed by a small decrease



This makes it possible to identify the number of eigenvalues to keep, and thus the dimensionality of the projections.

## CHOOSING d'

Eigenvalues measure the amount of distribution variance kept in the projection.

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Let us consider, for each *k* < *d*, the value

$$T_k = \frac{\sum_{i=1}^k \lambda_i^2}{\sum_{i=1}^n \lambda_i^2}$$

which provides a measure of the variance fraction associated to the *k* largest eigenvalues. When  $r_1 < \ldots < r_d$  are known, a certain amount *p* of variance can be kept by setting

 $d' = \operatorname*{argmin}_{i \in \{1, \dots, d\}} r_i > p$ 

## **PROBABILISTIC APPROACH TO PCA: IDEA**

Introduce a latent variable model to relate a *d*-dimensional observation vector to a corresponding d'-dimensional gaussian latent variable (with d' < d)

 $\mathbf{x} = \mathbf{W}\mathbf{z} + \boldsymbol{\mu} + \boldsymbol{\epsilon}$ 

where

- z is a d'-dimensional gaussian latent variable (the "projection" of x on a lower-dimensional subspace)
- W is a  $d \times d'$  matrix, relating the original space with the lower-dimensional subspace
- $\epsilon$  is a *d*-dimensional gaussian noise: noise covariance on different dimensions is assumed to be 0. Noise variance is assumed equal on all dimensions: hence  $p(\epsilon) = \mathcal{N}(\mathbf{0}, \sigma^2 \mathbf{I})$
- $\mu$  is the *d*-dimensional vector of the means

 $\epsilon$  and  $\mu$  are assumed independent.

## **GRAPHICAL MODEL**



1. 
$$\mathbf{z} \in \mathbb{R}^{d'}, \mathbf{x}, \epsilon \in \mathbb{R}^{d}, d' < d$$
  
2.  $p(\mathbf{z}) = \mathcal{N}(\mathbf{0}, \mathbf{I})$   
3.  $p(\epsilon) = \mathcal{N}(\mathbf{0}, \sigma^{2}\mathbf{I})$ , (isotropic gaussian noise)

## **GENERATIVE PROCESS**

This can be interpreted in terms of a generative process

1. sample the latent variable  $\mathbf{z} \in {\rm I\!R}^{d'}$  from

$$\boldsymbol{p}(\mathbf{z}) = \frac{1}{(2\pi)^{d'/2}} \boldsymbol{e}^{-\frac{||\mathbf{z}||^2}{2}}$$

2. linearly project onto  $\mathbb{R}^d$ 

$$\mathbf{y} = \mathbf{W}\mathbf{z} + \boldsymbol{\mu}$$

3. sample the noise component  $\boldsymbol{\epsilon} \in {\rm I\!R}^d$  from

$$\boldsymbol{p}(\boldsymbol{\epsilon}) = \frac{1}{(2\pi)^{d/2}} \boldsymbol{e}^{-\frac{||\boldsymbol{\epsilon}||^2}{2\sigma^2}}$$

4. add the noise component  $\epsilon$ 

$$\mathbf{x} = \mathbf{y} + \boldsymbol{\epsilon}$$

This results into  $p(\mathbf{x}|\mathbf{z}) = \mathcal{N}(\mathbf{W}\mathbf{z} + \boldsymbol{\mu}, \sigma^2 \mathbf{I})$ 

## **GENERATIVE PROCESS**



#### LATENT VARIABLE MODEL

The joint distribution is

$$p\left(\left[\begin{array}{c}\mathbf{z}\\\mathbf{x}\end{array}
ight]
ight)=\mathcal{N}(\boldsymbol{\mu}_{\mathbf{zx}},\boldsymbol{\Sigma})$$

By definition,

$$oldsymbol{\mu}_{\mathrm{zx}} = \left[ egin{array}{c} oldsymbol{\mu}_{\mathrm{z}} \ oldsymbol{\mu}_{\mathrm{x}} \end{array} 
ight]$$

- Since  $p(\mathbf{z}) = \mathcal{N}(\mathbf{0}, \mathbf{I})$ , then  $\mu_{\mathbf{z}} = 0$ .
- Since  $p(\mathbf{x}) = \mathbf{W}\mathbf{z} + \boldsymbol{\mu} + \boldsymbol{\epsilon}$ , then

$$\boldsymbol{\mu}_{\mathbf{x}} = \boldsymbol{E}[\mathbf{x}] = \boldsymbol{E}[\mathbf{W}\mathbf{z} + \boldsymbol{\mu} + \boldsymbol{\epsilon}] = \mathbf{W}\boldsymbol{E}[\mathbf{z}] + \boldsymbol{\mu} + \boldsymbol{E}[\boldsymbol{\epsilon}] = \boldsymbol{\mu}$$

Hence

$$\mu_{\mathrm{zx}} = \left[ egin{array}{c} \mathbf{0} \ \mathbf{\mu} \end{array} 
ight]$$

#### LATENT VARIABLE MODEL

For what concerns the distribution covariance

$$\mathbf{\Sigma} = \left[ egin{array}{cc} \mathbf{\Sigma}_{\mathrm{zz}} & \mathbf{\Sigma}_{\mathrm{zx}} \ \mathbf{\Sigma}_{\mathrm{zx}} & \mathbf{\Sigma}_{\mathrm{xx}} \end{array} 
ight]$$

where

$$\begin{split} \boldsymbol{\Sigma}_{zz} &= \boldsymbol{E}[\boldsymbol{z}\boldsymbol{z}^{\mathsf{T}}] = \boldsymbol{I} \\ \boldsymbol{\Sigma}_{zx} &= \boldsymbol{W}^{\mathsf{T}} \\ \boldsymbol{\Sigma}_{xx} &= \boldsymbol{W}\boldsymbol{W}^{\mathsf{T}} + \sigma^{2}\boldsymbol{I} \end{split}$$

#### LATENT VARIABLE MODEL

As a consequence, we get, for the joint distribution,

$$\boldsymbol{\mu}_{\mathrm{zx}} = \begin{bmatrix} \mathbf{0} \\ \boldsymbol{\mu} \end{bmatrix} \qquad \qquad \boldsymbol{\Sigma} = \begin{bmatrix} \mathbf{I} & \mathbf{W}^{\mathsf{T}} \\ \mathbf{W} & \mathbf{W}\mathbf{W}^{\mathsf{T}} + \sigma^{2}\mathbf{I} \end{bmatrix}$$

The marginal distribution of **x** is then  $p(\mathbf{x}) = \mathcal{N}(\boldsymbol{\mu}, \mathbf{W}\mathbf{W}^{\mathsf{T}} + \sigma^{2}\mathbf{I})$ 

The conditional distribution of z given x is  $p(z|x) = \mathcal{N}(\mu_{z|x}, \Sigma_{z|x})$  with

$$\begin{aligned} \boldsymbol{\mu}_{\mathbf{z}|\mathbf{x}} &= \mathbf{W}^{\mathsf{T}} (\mathbf{W} \mathbf{W}^{\mathsf{T}} + \sigma^{2} \mathbf{I})^{-1} (\mathbf{x} - \boldsymbol{\mu}) \\ \boldsymbol{\Sigma}_{\mathbf{z}|\mathbf{x}} &= \mathbf{I} - \mathbf{W}^{\mathsf{T}} (\mathbf{W} \mathbf{W}^{\mathsf{T}} + \sigma^{2} \mathbf{I})^{-1} \mathbf{W} = \sigma^{2} (\sigma^{2} \mathbf{I} + \mathbf{W}^{\mathsf{T}} \mathbf{W})^{-1} \end{aligned}$$

# MAXIMUM LIKELIHOOD FOR PCA

Setting  $\mathbf{C} = \mathbf{W}\mathbf{W}^{\mathsf{T}} + \sigma^{2}\mathbf{I}$ , the log-likelihood of the dataset in the model is

$$\log \mathbf{p}(\mathbf{X}|\mathbf{W}, \boldsymbol{\mu}, \sigma^2) = \sum_{i=1}^n \log \mathbf{p}(\mathbf{x}_i|\mathbf{W}, \boldsymbol{\mu}, \sigma^2)$$
$$= -\frac{nd}{2} \log(2\pi) - \frac{n}{2} \log |\mathbf{C}| - \frac{1}{2} \sum_{i=1}^n (\mathbf{x}_i - \boldsymbol{\mu}) \mathbf{C}^{-1} (\mathbf{x}_i - \boldsymbol{\mu})^T$$

Setting the derivative wrt  $\mu$  to zero results into

$$\boldsymbol{\mu} = \overline{\mathbf{x}} = \frac{1}{n} \sum_{i=1}^{n} \mathbf{x}_{i}$$

Maximization wrt W and  $\sigma^2$  is more complex: however, a closed form solution exists:

$$\mathbf{W} = \mathbf{U}_{d'} (\mathbf{L}_{d'} - \sigma^2 \mathbf{I})^{1/2}$$

where

- $\mathbf{U}_{d'}$  is the  $d \times d'$  matrix whose columns are the eigenvectors corresponding to the d' largest eigenvalues
- $\mathbf{L}_{d'}$  is the  $d' \times d'$  diagonal matrix of the largest eigenvalues

The columns of W are the principal components eigenvectors scaled by the variance  $\lambda_i - \sigma^2$ 

For what concerns maximization wrt  $\sigma^2$ , it results

$$\sigma^2 = rac{1}{oldsymbol{d} - oldsymbol{d}'} \sum_{i=oldsymbol{d}'+1}^{oldsymbol{d}} \lambda_i$$

since eigenvalues provide measures of the dataset variance along the corresponding eigenvector direction, this corresponds to the average variance along the discarded directions.

#### MAPPING POINTS TO SUBSPACE

The conditional distribution

$$\boldsymbol{p}(\mathbf{z}|\mathbf{x}) = \mathcal{N}(\mathbf{W}^{\mathsf{T}}(\mathbf{W}\mathbf{W}^{\mathsf{T}} + \sigma^{2}\mathbf{I})^{-1}(\mathbf{x} - \boldsymbol{\mu}), \sigma^{2}(\sigma^{2}\mathbf{I} + \mathbf{W}^{\mathsf{T}}\mathbf{W})^{-1})$$

can be applied.

In particular, the conditional expectation

 $\boldsymbol{E}[\mathbf{z}|\mathbf{x}] = \mathbf{W}^{\mathsf{T}}(\mathbf{W}\mathbf{W}^{\mathsf{T}} + \sigma^{2}\mathbf{I})^{-1}(\mathbf{x} - \boldsymbol{\mu})$ 

can be assumed as the latent space point corresponding to  $\mathbf{x}$ . The projection onto the d'-dimensional subspace can then be performed as

$$\mathbf{x}' = \mathbf{W}\mathbf{E}[\mathbf{z}|\mathbf{x}] + \boldsymbol{\mu} = \mathbf{W}\mathbf{W}^{\mathsf{T}}(\mathbf{W}\mathbf{W}^{\mathsf{T}} + \sigma^{2}\mathbf{I})^{-1}(\mathbf{x} - \boldsymbol{\mu}) + \boldsymbol{\mu}$$

Even if the log-likelihood has a closed form maximization, applying EM can sometimes be useful.

## **FACTOR ANALYSIS**

Noise components still gaussian and independent, but with different variance.



1. 
$$\mathbf{z} \in \mathbb{R}^d, \mathbf{x}, \epsilon \in \mathbb{R}^D, d \ll D$$
  
2.  $p(\mathbf{z}) = \mathcal{N}(\mathbf{0}, \mathbf{I})$   
3.  $p(\epsilon) = \mathcal{N}(\mathbf{0}, \Psi), \Psi$  diagonal (independent gaussian noise)

### **FACTOR ANALYSIS**

Model distribution are modified accordingly.

Joint distribution

$$p\left(\left[\begin{array}{c} \mathbf{z} \\ \mathbf{x} \end{array}\right]\right) = \mathcal{N}\left(\left[\begin{array}{c} \mathbf{0} \\ \mathbf{W} \end{array}\right], \left[\begin{array}{c} \mathbf{I} & \mathbf{W}^{\mathsf{T}} \\ \mathbf{\Lambda} & \mathbf{W}\mathbf{W}^{\mathsf{T}} + \mathbf{\Psi} \end{array}\right]\right)$$

Marginal distribution

$$p(\mathbf{x}) = \mathcal{N}(\boldsymbol{\mu}, \mathbf{W}\mathbf{W}^{\mathsf{T}} + \boldsymbol{\Psi})$$

Conditional distribution

The conditional distribution of z given x is now  $p(z|x) = \mathcal{N}(\mu_{z|x}, \Sigma_{z|x})$  with

$$\begin{split} \boldsymbol{\mu}_{\mathbf{z}|\mathbf{x}} &= \mathbf{W}^{\mathsf{T}} (\mathbf{W} \mathbf{W}^{\mathsf{T}} + \boldsymbol{\Psi})^{-1} (\mathbf{x} - \boldsymbol{\mu}) \\ \boldsymbol{\Sigma}_{\mathbf{z}|\mathbf{x}} &= \mathbf{I} - \mathbf{W}^{\mathsf{T}} (\mathbf{W} \mathbf{W}^{\mathsf{T}} + \boldsymbol{\Psi})^{-1} \mathbf{W} \end{split}$$

### MAXIMUM LIKELIHOOD FOR FA

The log-likelihood of the dataset in the model is now

$$\log p(\mathbf{X}|\mathbf{W}, \boldsymbol{\mu}, \boldsymbol{\Psi}) = \sum_{i=1}^{n} \log p(\mathbf{x}_i|\mathbf{W}, \boldsymbol{\mu}, \boldsymbol{\Psi})$$
$$= -\frac{nd}{2} \log(2\pi) - \frac{n}{2} \log |\mathbf{W}\mathbf{W}^{\mathsf{T}} + \boldsymbol{\Psi}| - \frac{1}{2} \sum_{i=1}^{n} (\mathbf{x}_n - \boldsymbol{\mu}) (\mathbf{W}\mathbf{W}^{\mathsf{T}} + \boldsymbol{\Psi})^{-1} (\mathbf{x}_i - \boldsymbol{\mu})^{\mathsf{T}}$$

Setting the derivative wrt  $\mu$  to zero results gain into

$$\boldsymbol{\mu} = \overline{\mathbf{x}} = rac{1}{n} \sum_{i=1}^{n} \mathbf{x}_i$$

Estimating parameters through log-likelihood maximization does not provide a closed form solution for W and  $\Psi$ . Iterative techniques such as EM must be applied.