

# **Machine Learning**

## **4 – Basics of Data Mining**

SS 2018

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1. Data preprocessing
  - Data format
  - Outlier detection
  - Missing values
2. Similarity measures

From now on, we mainly deal with continuous valued attributes.

Why do we need data preprocessing?

- Machine needs unique data format
- Outliers should be detected
- Missing values should be filled in

Aim:

Convert data of different formats and sources to a common format.

We will only **outline** the problem and not provide solutions.

Problem: Differing formats of attributes, e.g.,

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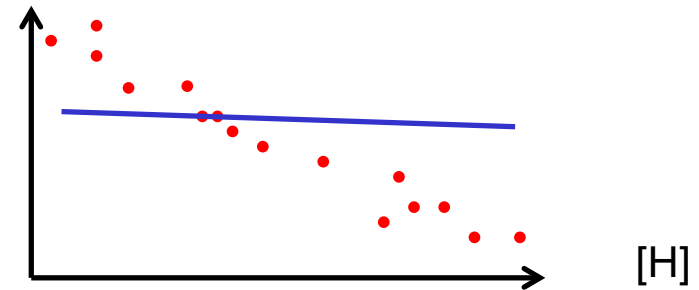
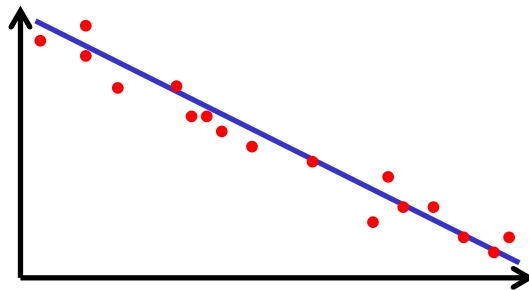
26 % - 0,26 - 0.26 - >1/4

- Unique identifiers such as the US social security number are often missing.
- Free formats
- Missing attributes
- Irrelevant additions
- Differing scales
- Systematic changes, e.g, change of name after marriage.
- Data formatted for different purpose. Example: Data records of persons, but you need records of households.

Major problems in data processing!

To date: **Only problem specific solutions.**

Outliers are rare – why is detection necessary ?



Trend of the regression line is spoiled by a single outlier

Outliers may have extreme values (often due to technical reasons) which can spoil statistics, in particular for small data sets.

Example:

Mean value  $\mu$  : 1,4 1,6 1,3 1,2 1,4 1,2 1,3  $\rightarrow \mu \approx 1,34$ .

Additional value: 7,3  $\rightarrow \mu \approx 2,09$ .

Some measures provide robustness against outliers even without explicit outlier detection.

Example: Replace mean by median  $m$  :

1,4   1,6   1,3   1,2   1,4   1,2   1,3    $\rightarrow \mu \approx 1,34.$

1,2   1,2   1,3   1,3   1,4   1,4   1,6    $\rightarrow m = 1,3.$

Outlier:

1,4   1,6   1,3   1,2   1,4   1,2   1,3   7,3    $\rightarrow \mu \approx 2,09.$

1,2   1,2   1,3   1,3   1,4   1,4   1,6   7,3    $\rightarrow m = 1,35.$

Causes of outliers:

1. Errors by measurement / technical errors
2. Unexpected “true” effect
3. Data with high variation – outliers are a natural part of the distribution

Effect related to 1:

Cut-off, e.g., limited range of measurement leads to high concentration of values at the boundaries.

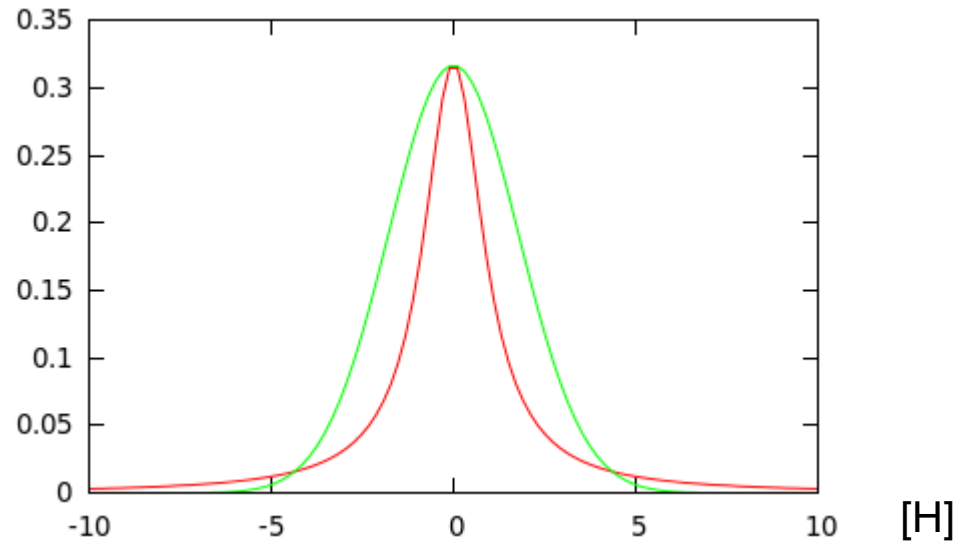
Problem 2 can be modeled by two or more overlaid distributions, e.g.,

$$P(x) = (1 - p) \cdot P_a(x) + p \cdot P_b(x), \quad p \ll 1.$$

Problem 3 can be modeled by a distribution with broad flanks, e.g.,

$$P(x) = 1 / (\pi(1+x^2)).$$





$$1 / (\pi(1+x^2))$$

$$1 / ((2\pi)^{1/2} \sigma) \cdot \exp(-1/2 (x/\sigma)^2), \quad \sigma=1.26$$

- Usually outliers are detected and removed.
- To consider a data point an outlier, we need to define what is *regular* !
- Most often a normal data distribution is assumed.
- For multivariate data, clustering algorithms can be applied and a normal distribution is assumed for each cluster.

Outliers of univariate distributions can be detected from  $z$ -values:

$$z_i = |x_i - \mu| / \sigma.$$

$z_i$  is a measure for the distance of  $x_i$  from the mean  $\mu$  in terms of the standard deviation  $\sigma$ .

Commonly, data with  $z_i > 3$  are considered outliers.

Improvement:

Outliers influence  $\mu$ , so use median instead. In this case, a threshold of 3,5 for outliers is used.

Detection of several outliers:

Idea: Iteratively remove outliers until z-tests finds no more.

1. Calculate mean  $\mu$  or median  $m$  and standard deviation  $\sigma$ .
2. Find data point  $x_{i^*}$  with largest z-value:  
$$i^* = \operatorname{argmax}_i z_i.$$
3. If  $x_{i^*}$  is an outlier, remove it from the data and goto 1.
4. Stop.

More efficient version:

Remove the  $k > 1$  outliers with largest z-values in each step.

Options:

- **Removal**: Simple, but loss of information.
- Don't remove outliers completely, but **weight** according to z-values.
- Remove and **fill up** gaps using the methods of the following section.

Why are missing values a major problem?

Example:

Data set of vectors from  $\mathcal{R}^{100}$ .

Probability that a value is missing is  $p = 2\%$ .

Thus probability that a vector is complete is  $(1 - p)^{100} = 13\%$ .

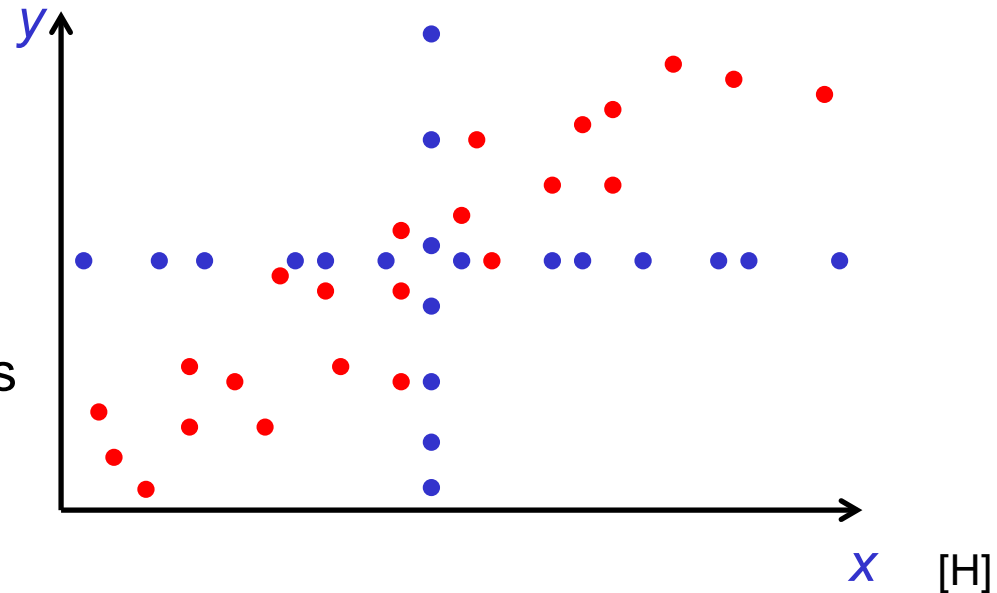
87% of all vectors are unusable unless we compute substitutes !

Problem: What can we do with missing values in a vector?

Example:

Data  $\{(x_i, y_i)\}$  show a clear trend.

How can missing values  
(i.e.,  $x$  or  $y$  are missing)  
be filled in ?



Idea 1:

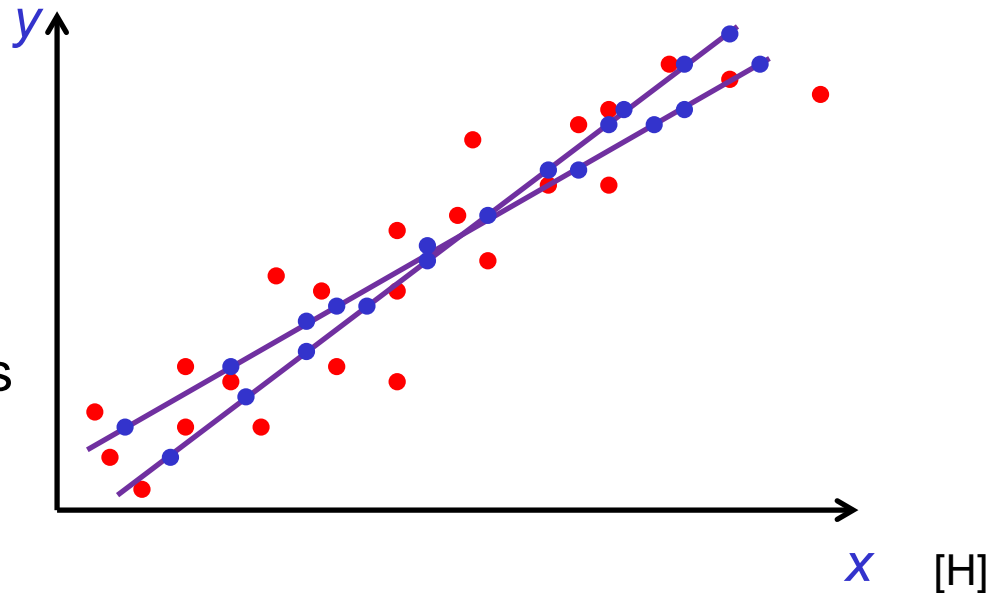
For data with only  $x$ , replace  $y$  by the mean or median  $\mu_y$   
from the rest of the data set (and vice versa) !

**Effect: Artifacts.**

Example:

Data  $\{(x_i, y_i)\}$  show a clear trend.

How can missing values (i.e.,  $x$  or  $y$  are missing) be filled in ?



Idea 2: Estimate a model to predict missing values.

Example: Linear regression

$$\begin{aligned} y' &= y_s \cdot x + y_0, & y_s &= C_{xy} / C_{xx}, & y_0 &= \mu_y - y_s \mu_x, \\ x' &= x_s \cdot y + x_0, & x_s &= C_{xy} / C_{yy}, & x_0 &= \mu_x - x_s \mu_y, \end{aligned}$$

where  $C_{xy} = \sum_{i=1..n} (x_i - \mu_x) \cdot (y_i - \mu_y)$  is the covariance of  $x$  and  $y$ .

**Effect: Artificial concentration of values along regression lines.**



## Why two different regression lines?

$$\begin{array}{lll} y' = y_s \cdot x + y_0, & y_s = C_{xy} / C_{xx}, & y_0 = \mu_y - y_s \mu_x. \\ x' = x_s \cdot y + x_0, & x_s = C_{xy} / C_{yy}, & x_0 = \mu_x - x_s \mu_y. \end{array}$$

Answer: Regression minimizes the mean square error of  $y$  depending on  $x$  (or alternatively  $x$  on  $y$ ).

So  $y'$  minimizes

$$\sum_{i=1..n} (y_i - (y_s \cdot x_i + y_0))^2 \rightarrow \min,$$

but  $x'$  minimizes

$$\sum_{i=1..n} (x_i - (x_s \cdot y_i + x_0))^2 \rightarrow \min,$$

with  $n = \# \text{ data}$ .

Do not confuse

$$C_{xy} = \sum_{i=1..n} (x_i - \mu_x) \cdot (y_i - \mu_y)$$

and Pearsons correlation coefficient

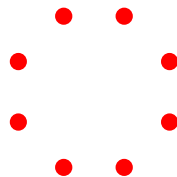
$$\rho_{xy} = C_{xy} / (\sigma_x \sigma_y),$$

with the standard deviation  $\sigma$ .

The correlation coefficient is a measure for linear dependence, taking values between  $-1$  (anti-correlation) and  $1$  (maximum correlation).

$0$  indicates independence.

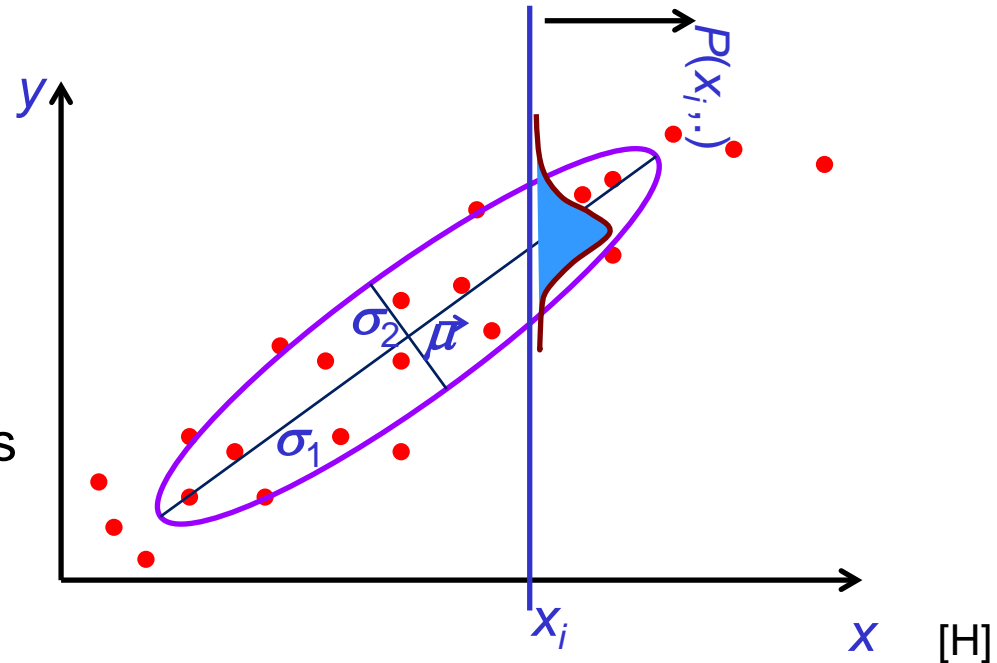
The absence of correlation  $\rho_{xy} = 0$  does not mean there is no structure in the data:



Example:

Data  $\{(x_i, y_i)\}$  show a clear trend.

How can missing values (i.e.,  $x$  or  $y$  are missing) be filled in ?



Idea 3:

Estimate the data distribution  $P(x, y)$ , e.g., assuming a normal distribution  $P(x, y) = N(\vec{\mu}, C)$ ,  $C$  = covariance matrix.

Once we have estimated  $\vec{\mu}$  and  $C$ , generate missing values  $x$  using a random number generator with probability distribution  $P(. , y)$ . and missing  $y$  from  $P(x, .)$ .

Problem when estimating  $P(x,y)$ :

We can use only the complete data vectors. The information of all data vectors with a missing value remains unused. This leads to

Idea 4: Estimate  $P$  using **Expectation Maximization (EM)**  
(Dempster, Laird, Rubin, 1977)

EM: Uses both complete and partial data vectors in an iterative procedure.

In the following, let

$x$  denote all **specified data** (complete and partial vectors),

$h$  all “**hidden**” (missing) values,

$\theta$  all **parameters** of the chosen distribution (such as mean and variances for a Gaussian).

The probability of the known values depends on the distribution (specified by  $\theta$ ) as  $P(x | \theta)$ .

The hidden values  $h$  are subject to the probability distribution (and thus to  $\theta$ ) as much as the known values  $x$ .

In addition, the hidden values  $h$  depend on  $x$ :  $P(h | x, \theta)$ .

The total distribution is thus

$$P(x, h | \theta) = P(h | x, \theta) \cdot P(x | \theta).$$

The *likelihood* of parameters  $\theta$  as a function of  $x$  and  $h$  is

$$L(\theta; x, h) = P(x, h | \theta).$$

For convenience, we consider the *log-likelihood* instead:

$$l(\theta) = \log L(\theta; x, h) = \log P(h | x, \theta) + \log P(x | \theta).$$

We want the parameters  $\theta^*$  that maximize the log-likelihood  $l(\theta)$ .

But  $l(\theta)$  depends on the hidden values  $h$ .

We get  $l(\theta) = \log P(x, h, \theta) = \log P(h | x, \theta) + \log P(x | \theta)$  by

- removing the hidden values  $h$  by “averaging out” to obtain an averaged  $\langle l(\theta) \rangle_h$ .
- To do so, we need the probability  $P(h | x, \theta)$ ,
- but this probability depends on the unknown  $\theta$ .
- So we need an estimate  $\theta_t$  for the real  $\theta$ .

The dilemma is solved by iteratively improving the estimate  $\theta_t$  for the real  $\theta$  (*M-step*) and averaging over  $h$  using the obtained  $\theta_t$  (*E-step*).  $\theta_t$  will converge to a local maximum  $\theta^*$  of  $l$  (hopefully close to  $\theta$ ).

Thus we maximize the averaged likelihood

$$Q(\theta, \theta_t) = \langle l(\theta) \rangle_h = \int P(h | x, \theta_t) \cdot \log P(h | x, \theta) dh + \log P(x | \theta)$$

So we have *traded* the  $h$ -dependence of  $l$  for a  $\theta_t$ -dependence of  $Q$ .

Procedure:

1. Choose a function to approximate  $P(x, y \mid \theta)$  with parameters  $\theta$ .
2. Choose start values  $\theta_t$ .
3. Initialize step counter  $t = 0$ .

4. **E-step**: Calculate the integral of

$$Q(\theta, \theta_t) = \int P(h \mid x, \theta_t) \cdot \log P(h \mid x, \theta) dh + \log P(x \mid \theta)$$

to obtain the function  $Q$  depending on  $\theta$  and  $\theta_t$ .

5. **M-step**: Maximize  $Q$  with respect to  $\theta$ :

$$\theta_{t+1} = \arg \max_{\theta} Q(\theta, \theta_t).$$

6.  $t++$ .

7. If  $Q(\theta, \theta_t)$  does not meet the convergence condition goto 4.

Steps 4 and 5 may include heavy numerics.

# EM-algorithm: Proof of convergence

Show that the data log-likelihood  $\log P(x|\theta)$  increases with each EM-step.

Let's use  $\mathcal{P}_t$  for short for  $P(h|x, \theta_t)$  and  $\mathcal{P}$  for  $P(h|x, \theta)$ , so we get

$$Q(\theta, \theta_t) = \int P(h|x, \theta_t) \cdot \log P(h|x, \theta) dh + \log P(x|\theta) = \int \mathcal{P}_t \cdot \log \mathcal{P} dh + \log P(x|\theta).$$

Change  $\Delta_t$  of  $\log P(x|\theta) = Q(\theta, \theta_t) - \int \mathcal{P}_t \cdot \log \mathcal{P} dh$  under one EM-step:

$$\Delta_t := \log P(x|\theta_{t+1}) - \log P(x|\theta_t)$$

$$= Q(\theta_{t+1}, \theta_t) - \int \mathcal{P}_t \cdot \log \mathcal{P}_{t+1} dh - Q(\theta_t, \theta_t) + \int \mathcal{P}_t \cdot \log \mathcal{P}_t dh$$

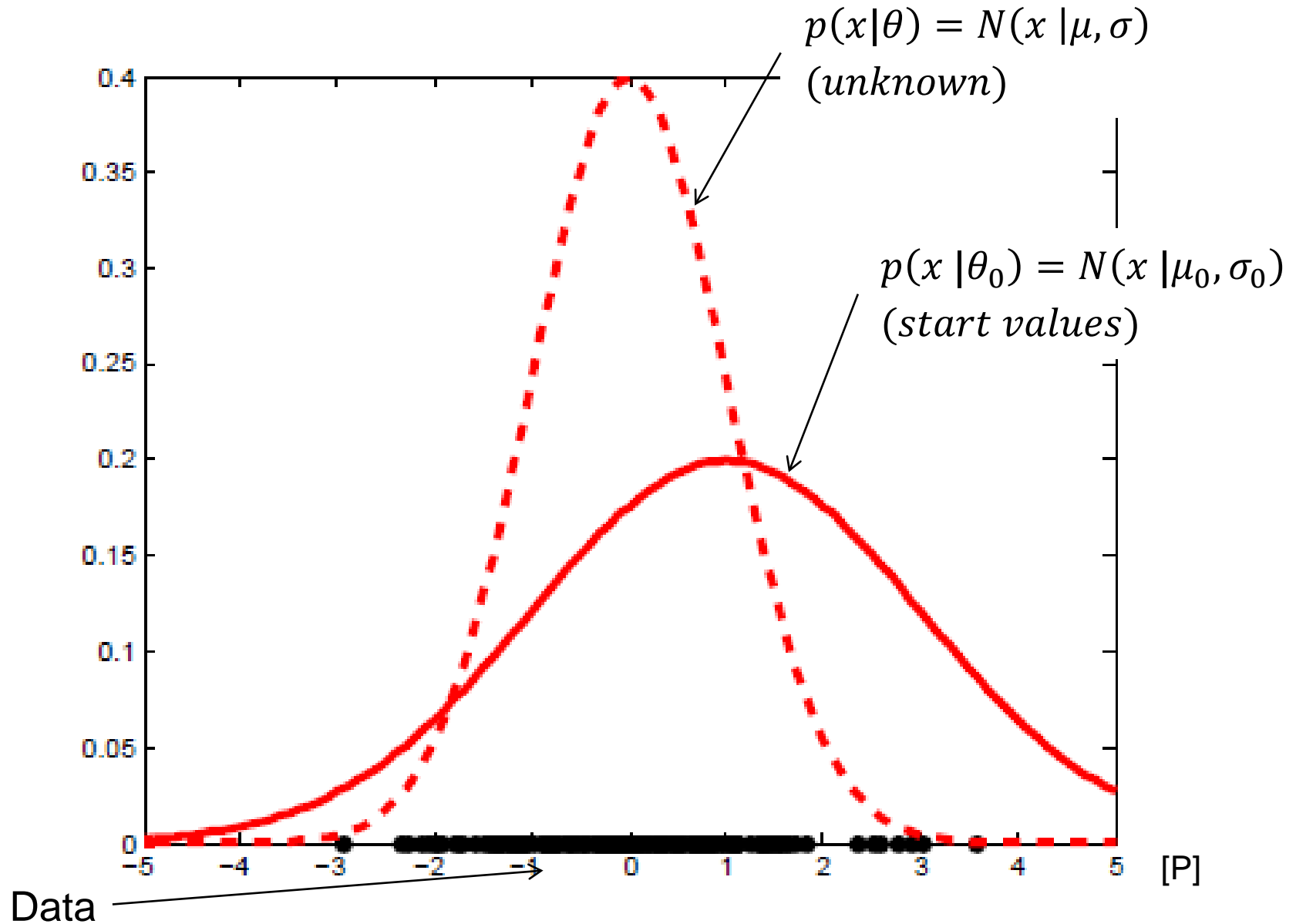
$Q(\theta_{t+1}, \theta_t) - Q(\theta_t, \theta_t) =: Q_t \geq 0$  holds for all  $t$  by definition of the M-step.

$$\Delta_t = Q_t + \int \mathcal{P}_t \cdot \log \mathcal{P}_t dh - \int \mathcal{P}_t \cdot \log \mathcal{P}_{t+1} dh = Q_t + \int \mathcal{P}_t \cdot \log(\mathcal{P}_t / \mathcal{P}_{t+1}) dh \geq 0$$

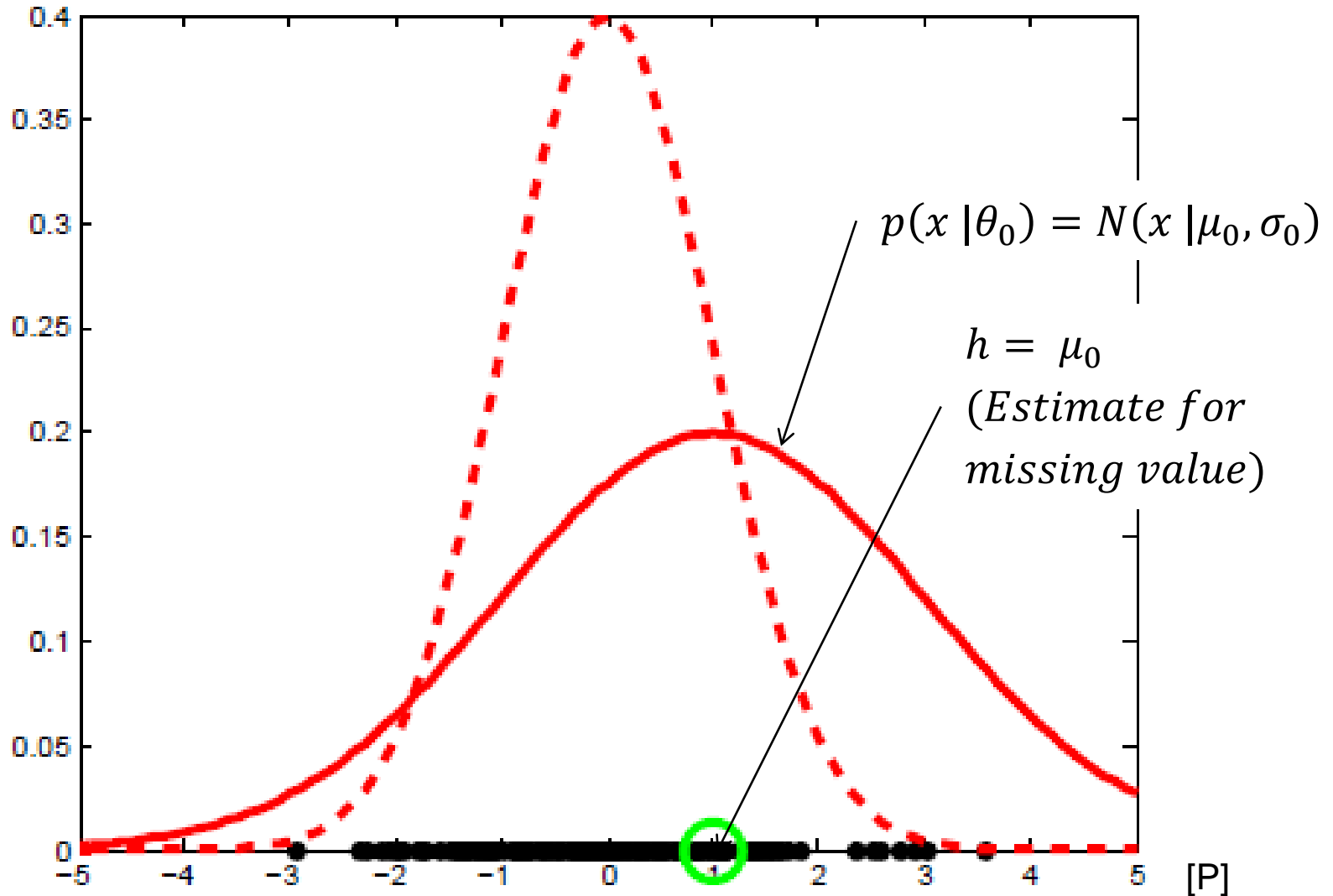
holds for all  $t$  because the integral is the Kullback-Leibler divergence of  $\mathcal{P}_t$  and  $\mathcal{P}_{t+1}$ , which is always non-negative (Gibbs inequality).



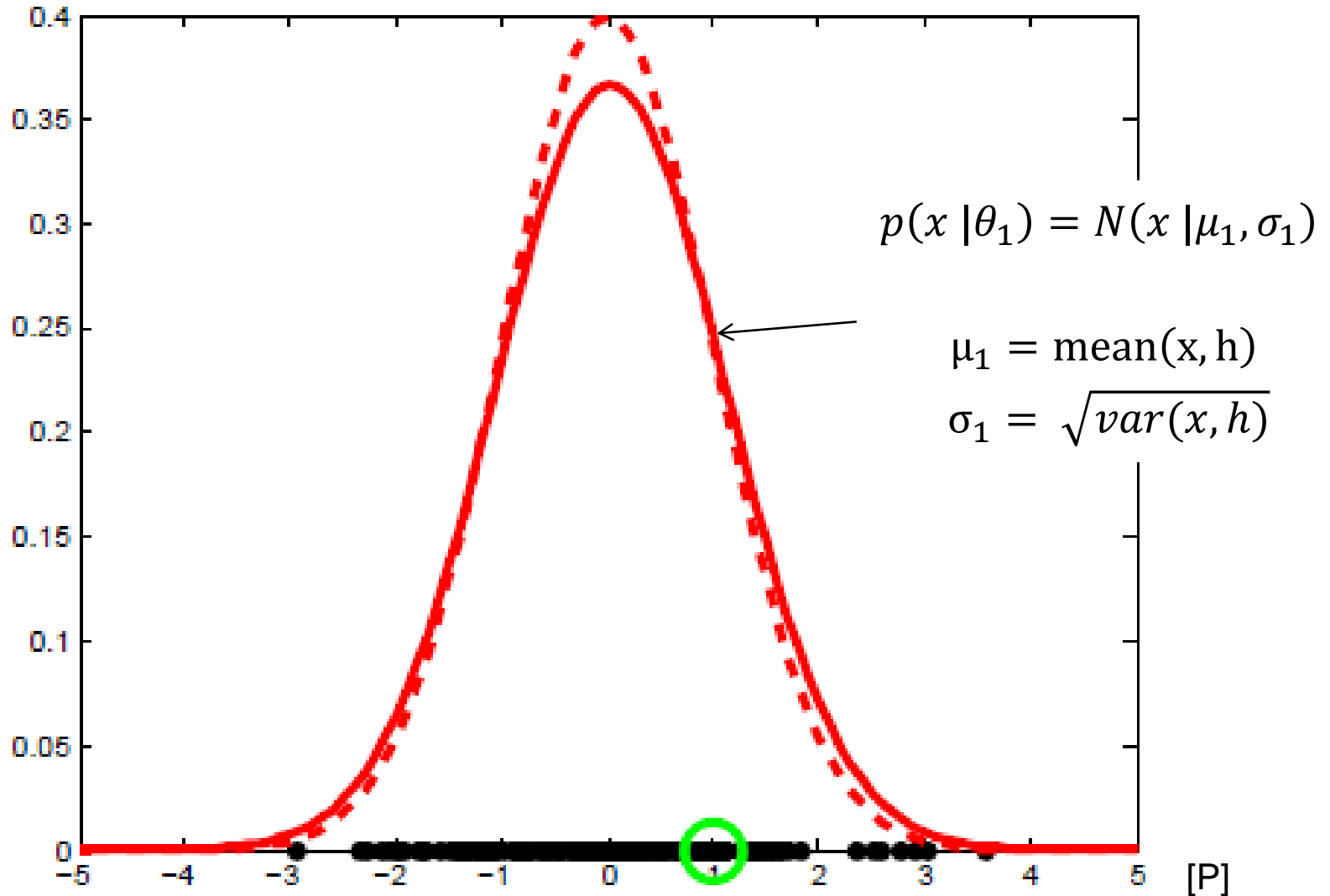
# EM algorithm – 1D-Example



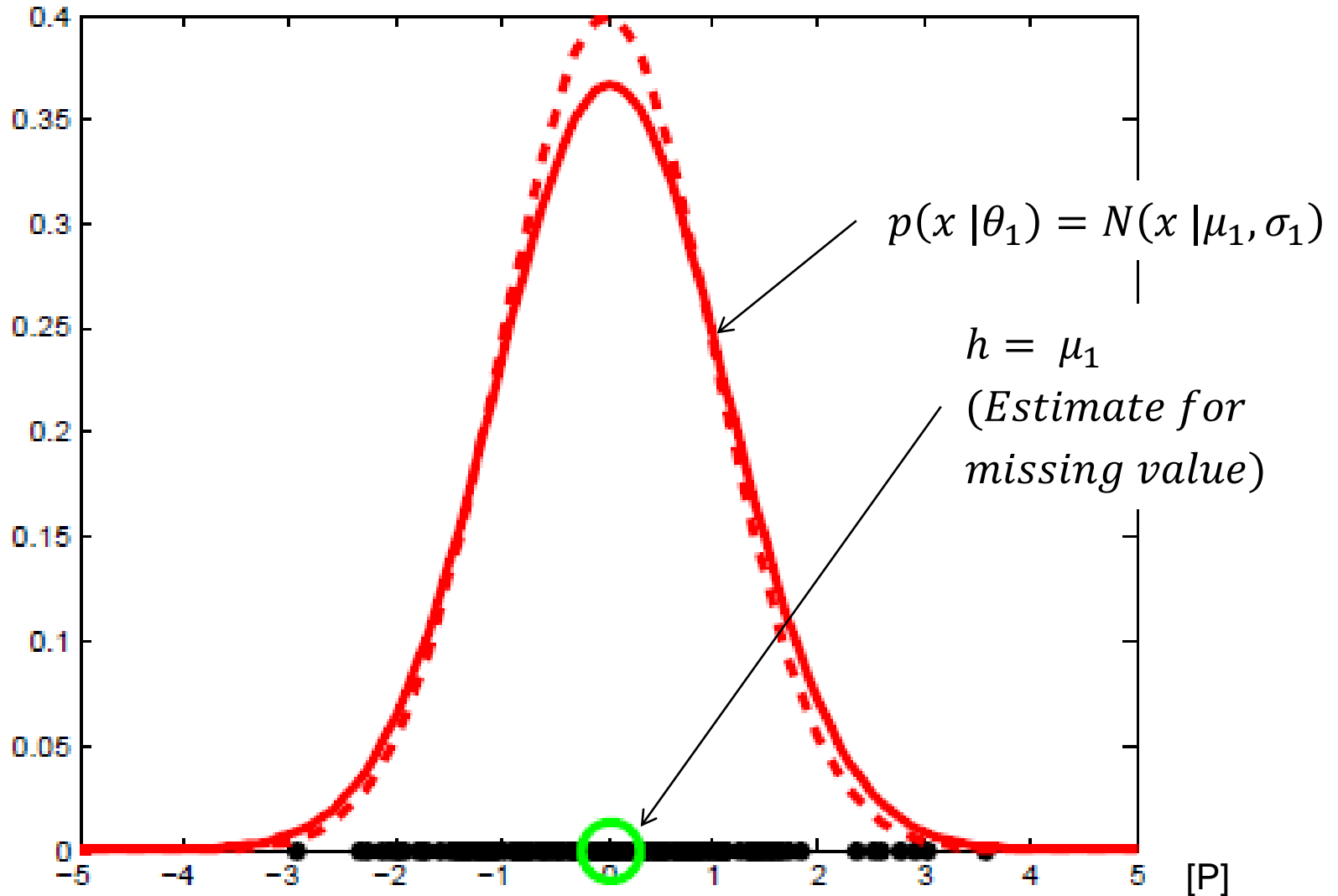
# EM algorithm – Expectation step



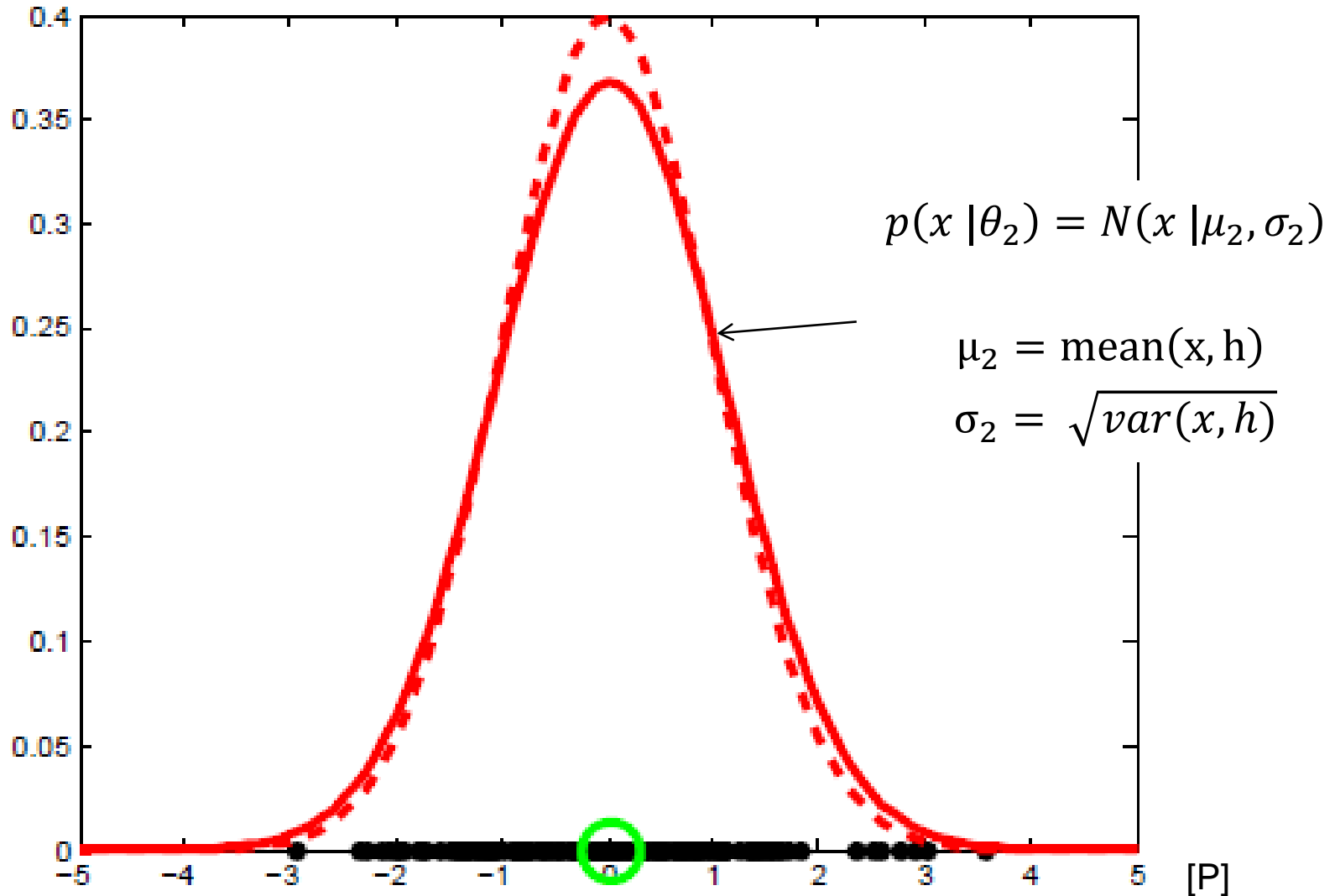
# EM algorithm – Maximisation step



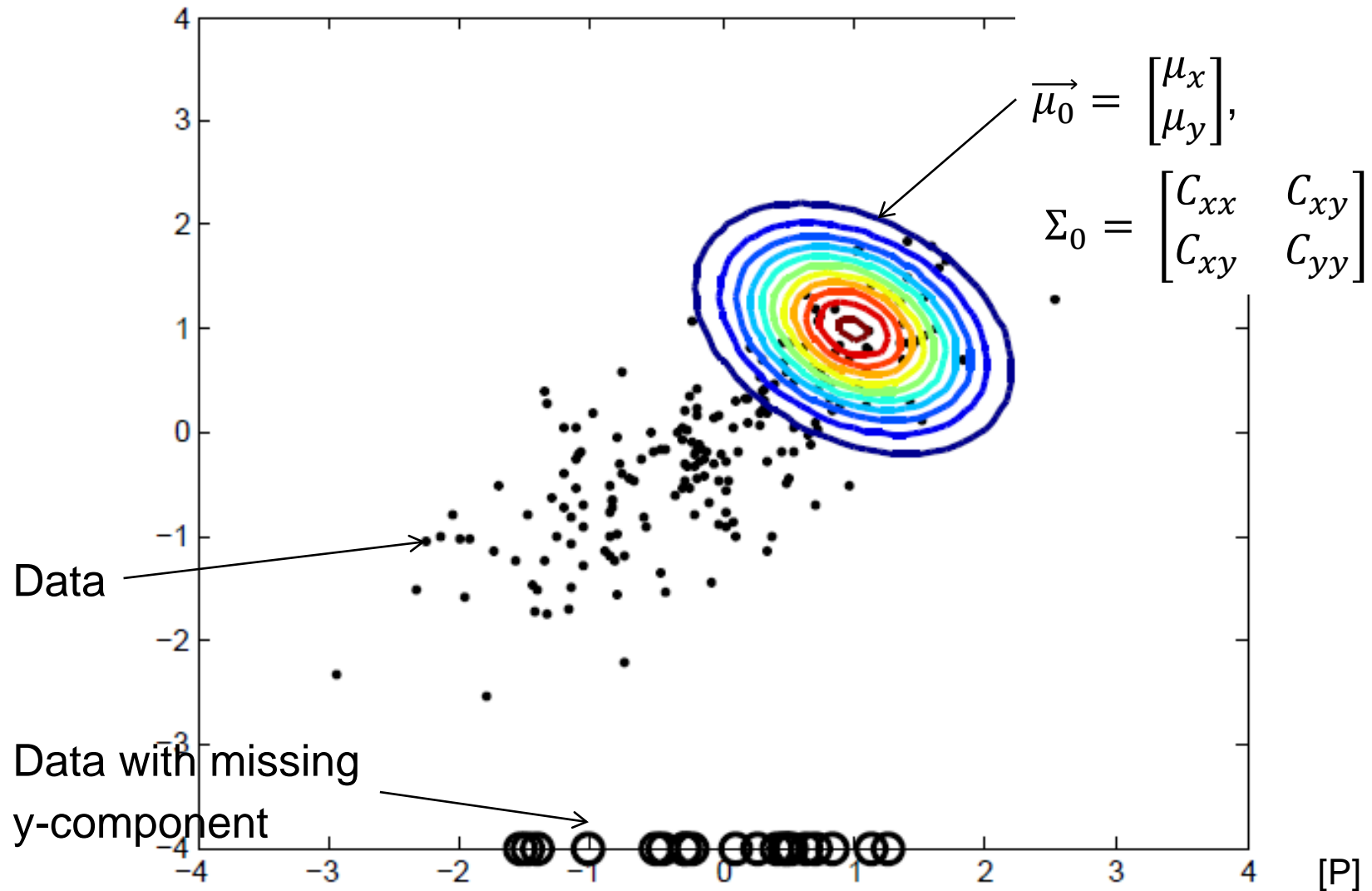
## EM algorithm – Expectation step 2



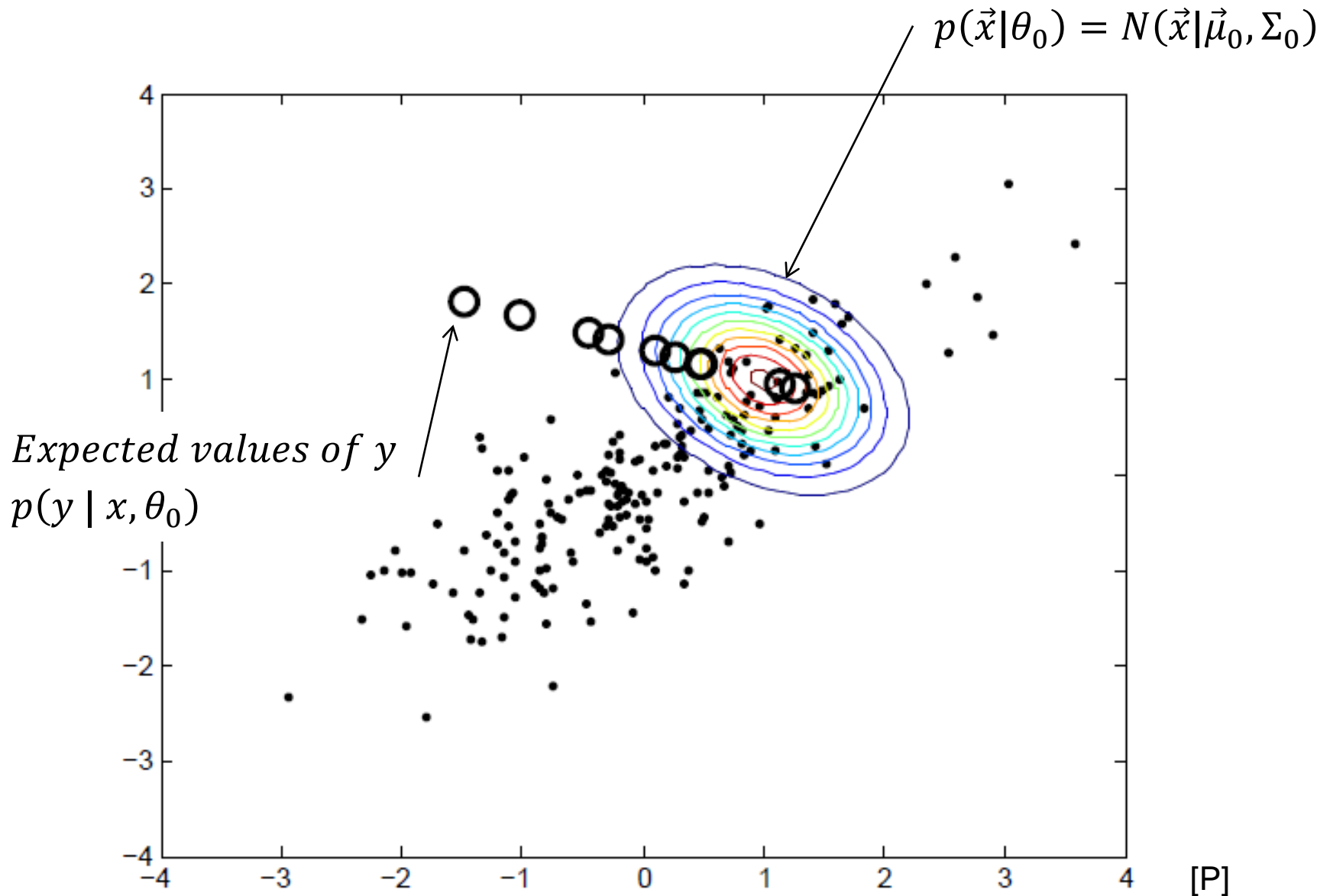
## EM algorithm – Maximisation step 2



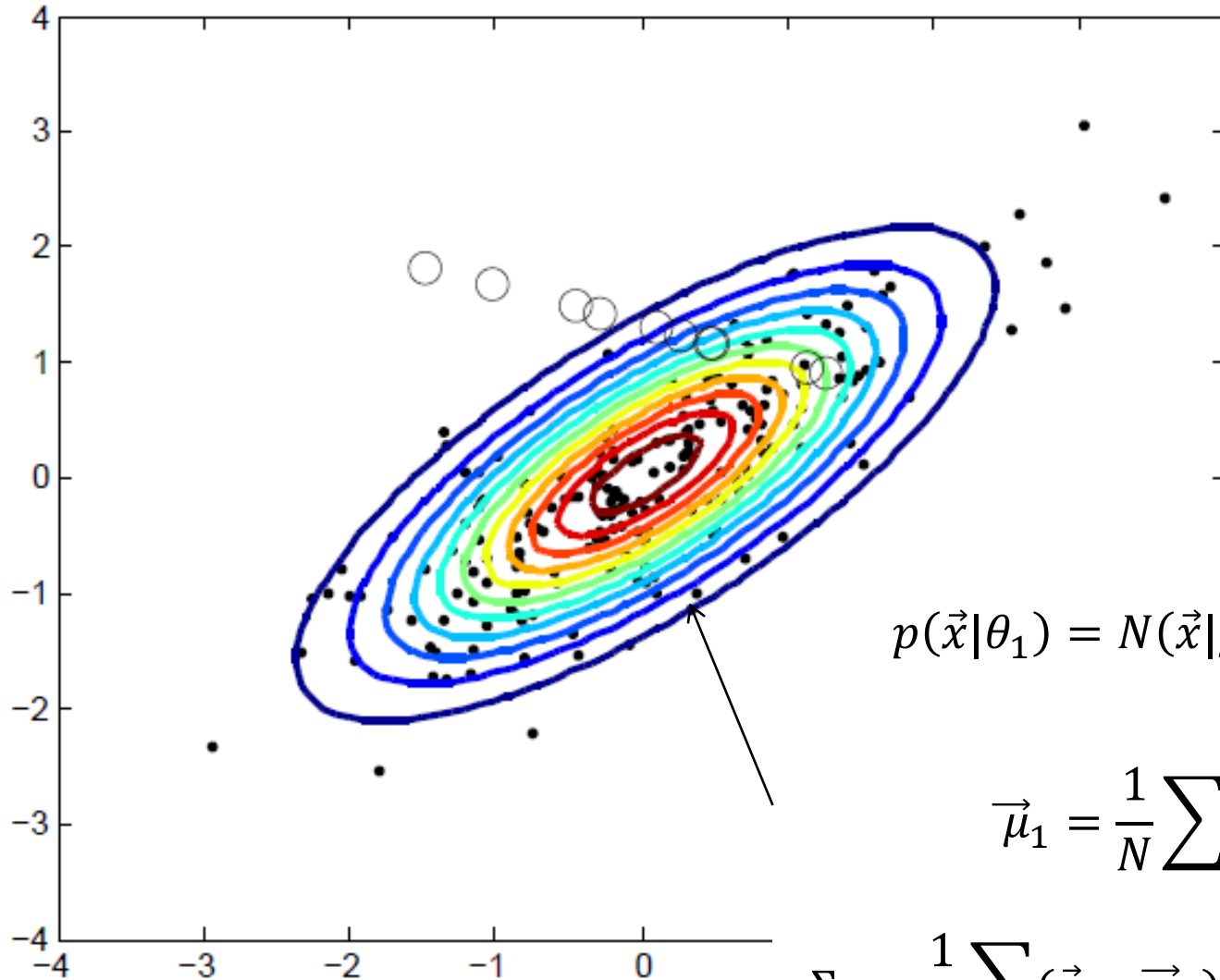
$$p(\vec{x}|\theta_0) = N(\vec{x}|\vec{\mu}_0, \Sigma_0)$$



# EM algorithm – Expectation Step



# EM algorithm – Maximisation Step



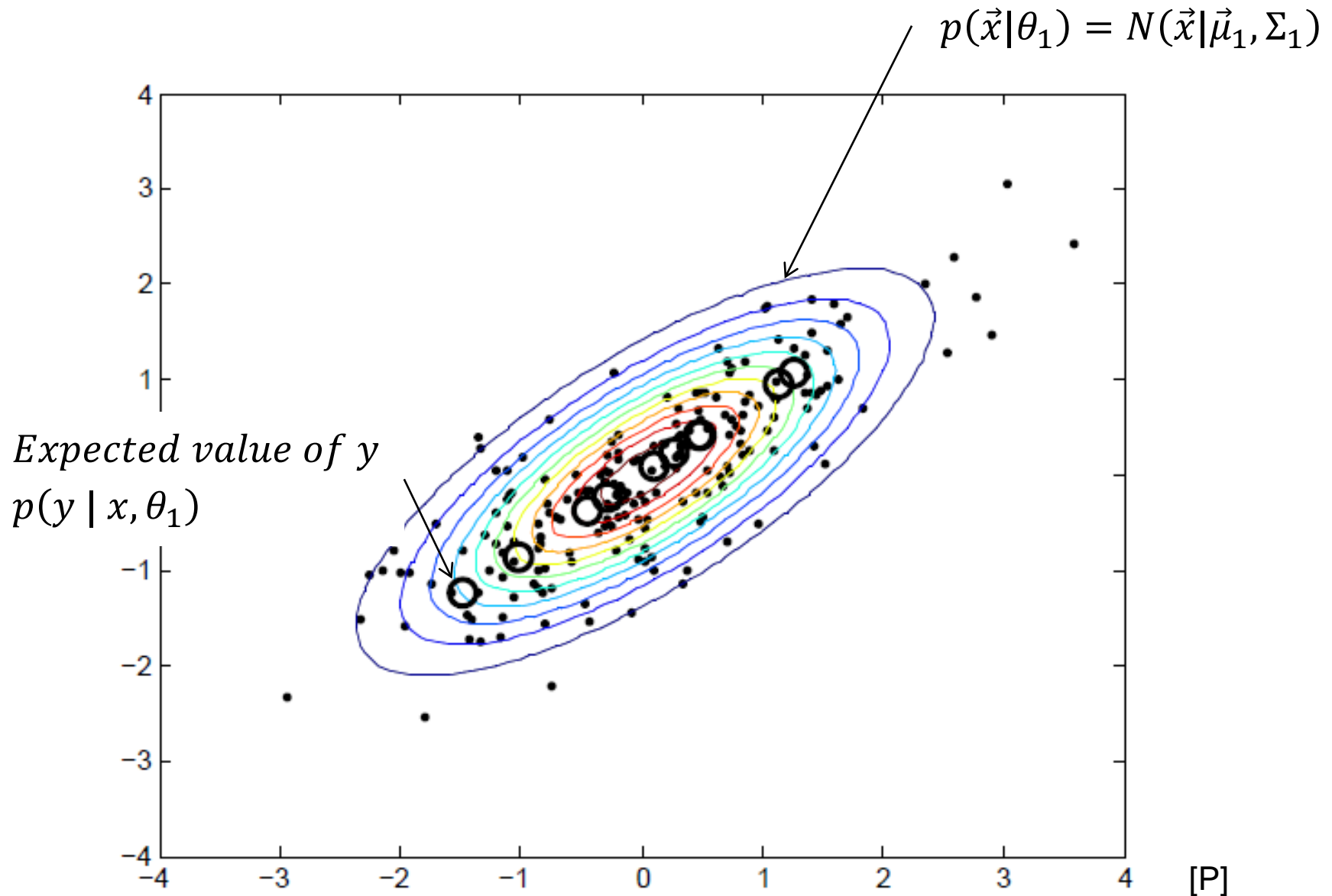
$$p(\vec{x}|\theta_1) = N(\vec{x}|\vec{\mu}_1, \Sigma_1)$$

$$\vec{\mu}_1 = \frac{1}{N} \sum \vec{x}_i$$

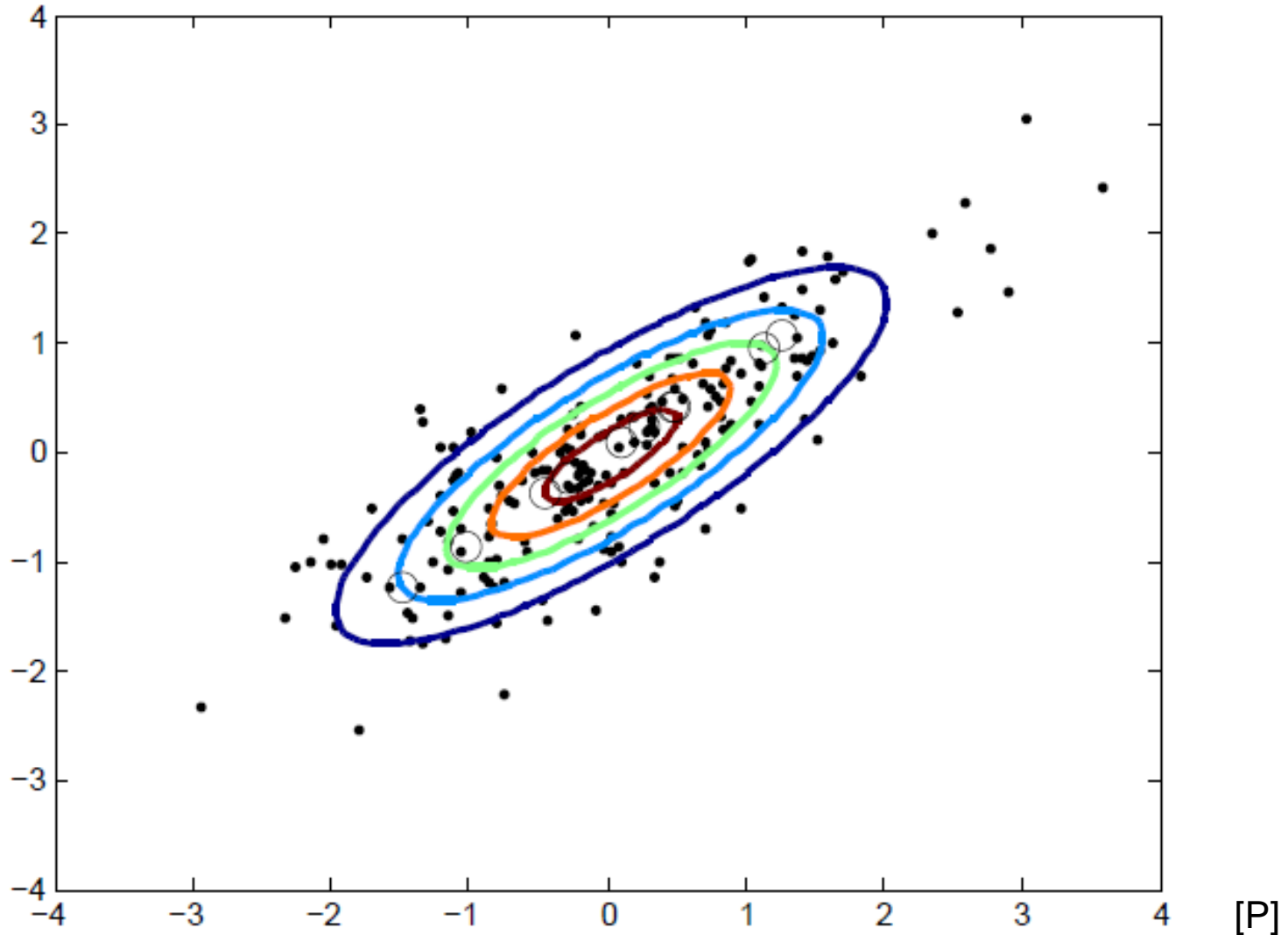
$$\Sigma_1 = \frac{1}{N} \sum (\vec{x}_i - \vec{\mu}_1)(\vec{x}_i - \vec{\mu}_1)^T \quad [\text{P}]$$



# EM algorithm – Expectation Step



# EM algorithm – Maximisation Step



- Missing information can not be regained.
- We can only „invent“ missing values such that they do not contradict the existing ones. For this we use models.
- At best, we do not destroy information by this procedure (but usually we do).
- Then why did we care about missing value substitution?

Answer:      To make the existing values technically usable!

- Data format is major issue for applications.
- Outlier detection requires some definition of what is regular.
- Missing values can be a major problem when spread among large data records, rendering most records unusable.
- Missing value substitution by mean of other data yields artifacts.
- Fitting a model to the data by regression is better, but does not make use of partial data.
- EM-algorithm solves this problem but suffers from problems of local search.

# Similarity measures

What we really want:

**Relations** between data on the **semantic level**, in particular, **similarity / dissimilarity**.

What is accessible to a machine:

Numerical measures, in particular, *distance functions* (*metrics*).

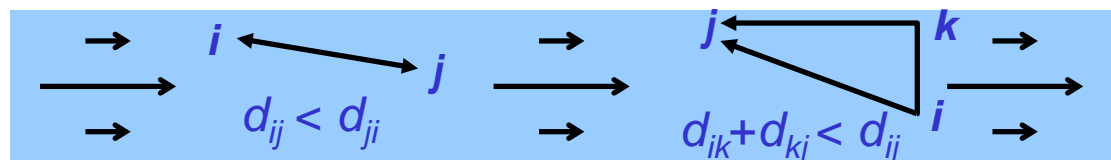
A **distance function** or **metric**  $d$  must obey the following conditions, which are reasonable for **geometric** distances ( $i$  to  $j$  are locations):

1. Symmetry:  $d_{ij} = d_{ji}$  (from  $i$  to  $j$  is as far as from  $j$  to  $i$ ).
2. Coincidence axiom:  $d_{ij} = 0 \Leftrightarrow i = j$  (identity of indiscernibles).
3. Triangle equation:  $d_{ik} + d_{kj} \geq d_{ij}$  (way over  $k$  is no shorter than direct path from  $i$  to  $j$ ).

Note the axioms imply  $d_{ij} \geq 0$  (non-negativity).

**Question:** Think of a “distance related” quantity that is not metric!

$d_{ij}$  is fuel consumption of a vessel going from  $i$  to  $j$  in a river.



Remark:

In mathematics, the term *distance function* is used only when the axioms are fulfilled.

In ML, *distance function* is often used like *dissimilarity function* and may be applied to quantities that do not match the axioms.

To make it crystal clear you mean a distance function fulfilling the axioms, use the term *metric*.



For a data set  $\{\vec{x}_1 \dots \vec{x}_n\}$ , all information about distances is assembled in the distance matrix

$$D = \begin{bmatrix} d_{11} & d_{12} & \dots & d_{1n} \\ \vdots & & & \vdots \\ d_{n1} & d_{n2} & \dots & d_{nn} \end{bmatrix}$$

where  $d_{ij}$  denotes the distance between  $\vec{x}_i$  and  $\vec{x}_j$ .

Note  $d_{ii} = 0$  and  $d_{ij} = d_{ji}$ .

Distance calculation is motivated from **geometric** distances.

But in ML, distances are more broadly used to express **similarity**.

Similarities of data may be represented by a matrix as well as distances.

When similarities are not computed from features of the data, but assigned explicitly from other sources (e.g., human insight), similarities may become particularly “non-geometric”.

Example: *x likes y* on a scale 1...10 :

No attributes specified  
for the persons,  
only distances!



likes	Luke	Leia	Han
Luke	8	9	6
Leia	7	8	10
Han	6	9	16

Some common distances for  $\vec{x}, \vec{y} \in \mathbb{R}^L$ :

Euclidean distance:

$$d(\vec{x}, \vec{y}) = \|\vec{x} - \vec{y}\| = \left( \sum_{i=1 \dots L} (x_i - y_i)^2 \right)^{\frac{1}{2}}$$

- Simple and frequently used measure.
- No individual weighting of components.

Example: Broomstick production

Produced broomstick of dimensions  $x = (\text{length}, \text{thickness})^T$  is compared to prototype values  $y = (160\text{cm}, 3\text{cm})^T$ .

Problem:  $\|(161\text{cm}, 3\text{cm})^T - y\| = \|(160\text{cm}, 4\text{cm})^T - y\|$ .

Idea: Weight dimensions according to variation from the mean.

Normalized euclidean distance (also: Pearson or  $\chi^2$ -distance):

$$d(\vec{x}, \vec{y}) = \left( \sum_{i=1 \dots L} (x_i - y_i)^2 / \sigma_i \right)^{1/2}$$

with standard deviations  $\sigma_i$ .

**Question:** How does the broomstick example motivate the Pearson distance ?

- At first glance, the example suggests weighting according to mean, not deviation.
- But with smaller absolute value, the precision of production usually increases.

Problem: Correlated vector components

Example:

$$\vec{x}(t) = (x_1(t), x_2(t), \dots, x_L(t))^T \quad \text{with}$$
$$x_1(t) = x_2(t) = \dots = x_{L-1}(t), \quad \sigma_1 = \dots = \sigma_L = 1.$$

Pearson distance:

$$d(\vec{x}(t_1), \vec{x}(t_2)) = ((L-1) \cdot (x_1(t_1) - x_1(t_2))^2 + 1 \cdot (x_L(t_1) - x_L(t_2))^2)^{1/2}$$

→ The correlated components are over-weighted.

Idea: Scaling of distances using the covariance matrix  $C$ .

$$d(\vec{x}, \vec{y}) = ((\vec{x} - \vec{y})^T C^{-1} (\vec{x} - \vec{y}))^{1/2}$$

Properties:

- Scale and translation invariant.
- If  $C$  is unit matrix: Euclidean distance.
- Points of equal Mahalanobis distance to a center form an ellipsoid.
- Scaling might destroy structure within the data.

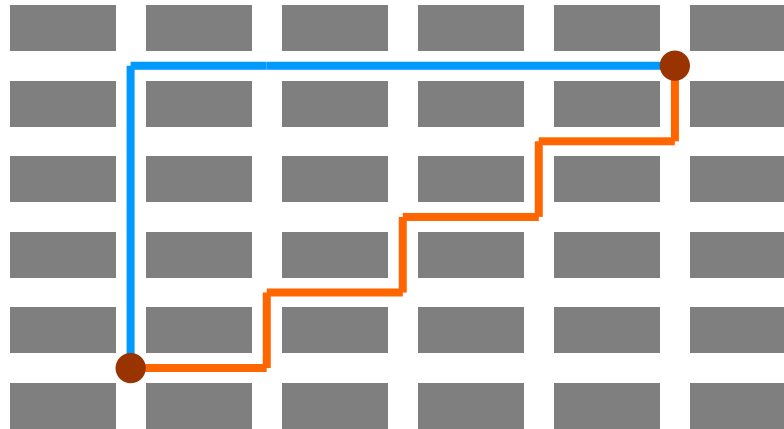
Interpretation: Diagonalize  $C$

$$C = \begin{bmatrix} \lambda_1 & & & 0 \\ & \lambda_2 & & \\ & & \ddots & \\ 0 & & & \lambda_L \end{bmatrix}$$

with eigenvalues  $\lambda_i \rightarrow$  scaling factors are  $(\lambda_i)^{-1/2}$  !

Also: *Manhattan distance*

$$d(\vec{x}, \vec{y}) = \sum_{i=1 \dots L} |x_i - y_i|$$



[H]

Remark:

The *Hamming distance* (= number of positions where two strings of equal length differ) is equal to the Manhattan distance for binary strings.

Also: *Maximum distance, chessboard distance*

$$d(\vec{x}, \vec{y}) = \max_{i=1 \dots L} |x_i - y_i|$$

Minimum number of moves a king needs between two positions on a chessboard.



Generalization:

$$d(\vec{x}, \vec{y}) = \left( \sum_{i=1 \dots L} |x_i - y_i|^p \right)^{1/p}$$

Special cases:

$$p = 1: \quad d(\vec{x}, \vec{y}) = \sum_{i=1 \dots L} |x_i - y_i| \quad (\text{city block})$$

$$p = 2: \quad d(\vec{x}, \vec{y}) = \left( \sum_{i=1 \dots L} (x_i - y_i)^2 \right)^{1/2} \quad (\text{euclidean})$$

$$p \rightarrow \infty: \quad d(\vec{x}, \vec{y}) = \max_{i=1 \dots L} |x_i - y_i| \quad (\text{maximum})$$

Unit circles (schematic):



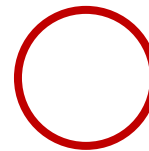
$p$  small



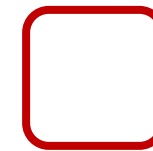
$p = 1/2$



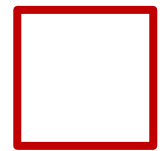
$p = 1$



$p = 2$



$p$  big



$p \rightarrow \infty$

[H]

So far, all distance measures relied on the topology of an  $\mathbb{R}^n$ .

There are data with other topologies, e.g., angular attributes (topology of a circle):

$$||10^\circ - 30^\circ|| = 20^\circ$$

$$||0^\circ - 359^\circ|| = 359^\circ$$

Solution: *Embedding* complex topologies into an  $\mathbb{R}^n$ .

Mapping of nominal attribute values to real values:

*(low, medium, high)*  $\rightarrow$  (1, 2, 3) makes sense, but

*(stone, wood, metal)*  $\rightarrow$  (1, 2, 3) implies an order that is not there.

Solution: *(stone, wood, metal)*  $\rightarrow$   $((1,0,0)^T, (0,1,0)^T, (0,0,1)^T)$

Problem:

For a large number  $n$  of attribute values, dimensionality becomes too high.

Solution: Choose normalized random vectors  $v_i \in \mathbb{R}^m, i = 1 \dots n$ ,  $1 \ll m \ll n$  instead. Vectors drawn at random from a space of high dimension tend to be close to orthogonal (why?).

For binary attributes (e.g., *small / big*, *yes / no*) use the Jaccard index  $J$  as a similarity measure which is defined for sets  $A$  and  $B$ :

$$\begin{aligned} J(A, B) &= (\# \text{ common elements}) / (\# \text{ all elements}) \\ &= |A \cap B| / |A \cup B| \end{aligned}$$

Example:

$$J(\{a, c, d, e\}, \{b, c, e, f\}) = 2 / 6$$

Jaccard distance function:

$$J_d(A, B) = 1 - J(A, B)$$

Application in text mining:

Similarity of strings can be calculated by cutting strings into tokens and using  $J$  on the token sets.

- [M] Online material available at [www.cs.cmu.edu/~tom/mlbook.html](http://www.cs.cmu.edu/~tom/mlbook.html) for the textbook: Tom M. Mitchell: *Machine Learning*, McGraw-Hill
- [H] Gunther Heidemann, 2012.
- [P] Michael Pardowitz, 2014.