## MVA – Probabilistic Graphical Models

## Homework 2

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## K-means and the EM algorithm 1

**Question 1.** We consider a mixture model of K components for a dataset  $(X_i)$  where  $X_i \in \mathbb{R}^d$ . We denote  $Z_i \in \{1, \dots, K\}$  the latent hidden label. Each component occurs with probability  $p_k = \mathbb{P}(Z_i = k)$  and is distributed as

$$X_i \sim \mathcal{N}(\mu_k, D_k)$$

i.e.  $p(x|k) = \frac{1}{((2\pi)^d |D_k|)^{1/2}} \exp(-\frac{1}{2}(x - \mu_k)^T D_k^{-1}(x - \mu_k)).$ The data log-likelihood under parameters  $\Theta = ((p_1, \mu_1, D_1), \dots, (p_K, \mu_K, D_K))$  is

$$\mathcal{L}(X_1, \dots, X_n; \Theta) = \sum_{i=1}^n \log \left( \sum_{k=1}^K p_k p(X_i | k; \mu_k, D_k) \right)$$

We seek to compute the MLE

$$\widehat{\Theta} \in \operatorname*{argmax}_{\Theta} \mathcal{L}(X_1, \dots, X_n; \Theta)$$

This optimization problem is intractable when using straightforward methods. The EM algorithm goes as follows:

Expectation Compute the posterior probability of the latent variables  $Z_i$ :

$$q_{k,i}^{(t)} = p(Z_i = k|X_i; \Theta^{(t)}) = \frac{p_k^{(t)} p(X_i|Z_i = k; \Theta^{(t)})}{\sum_{\ell=1}^K p_\ell^{(t)} p(X_i|Z_i = \ell; \Theta^{(t)})}$$
(1)

and denote  $w_k^{(t)} = \sum_{i=1}^n q_{k,i}^{(t)}$  - we then have  $\sum_{k=1}^K w_k^{(t)} = n$ .

• Maximization Update the parameters  $\Theta^{(t)}$  by maximizing the lower bound objective:

$$\max_{\Theta} \mathcal{J}(q^{(t)}, \Theta) = \sum_{i=1}^{n} \left( \sum_{k=1}^{K} q_{k,i}^{(t)} \left( \log p_k - \frac{d}{2} \log(2\pi) - \frac{1}{2} \log|D_k| - \frac{1}{2} (X_i - \mu_k)^T D_k^{-1} (X_i - \mu_k) \right) \right)$$
(2)

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subject to  $\sum_{k=1}^{K} p_k = 1$  (associated to a multiplier  $\nu$ ). The KKT conditions give the null gradient condition:

$$\frac{1}{p_k} w_k^{(t)} - \nu = 0 (3a)$$

$$-\sum_{i=1}^{n} q_{k,i}^{(t)} D_k^{-1} (\mu_k - X_i) = 0$$
(3b)

$$-\frac{1}{2}\sum_{i=1}^{n} q_{k,i}^{(t)} (D_k^{-1} - D_k^{-2} \operatorname{diag}(X_i - \mu_k)^2) = 0$$
(3c)

Which leads to the updates:

$$p_k = \frac{1}{n} w_k^{(t)} \tag{4a}$$

$$\mu_k = \frac{1}{w_k^{(t)}} \sum_{i=1}^n q_{k,i}^{(t)} X_i \tag{4b}$$

$$D_k = \frac{1}{w_k^{(t)}} \sum_{i=1}^n q_{k,i}^{(t)} \operatorname{diag}(X_i - \mu_k)^2$$
 (4c)

**Question 2.** The main advantage of this "reduced" covariance mixture model is that it is more sparse: it uses far fewer parameters (K(2d+1)) than the its full counterpart, which has K(1+d+d(d+1)/2) parameters. For datasets with relatively independent features (conditionally on the latent class), this can give performance very close to the full covariance while having a smaller, simpler model (meaning better AIC or BIC scores).

**Question 3.** Figure 1 compares the obtained latent class centroids and confidence ellipsoids (where applicable) for the diagonal and full covariance mixture models and K-means, on the Iris dataset, for a small number of classes K = 3 (the actual number of classes in the data). Figures 2 and 3 represent the same for K = 2, 4 classes.

**Question 4.** I created synthetic data made up of an ellipsis and annulus: Figure 4 shows how the mixture models and K-means compare.

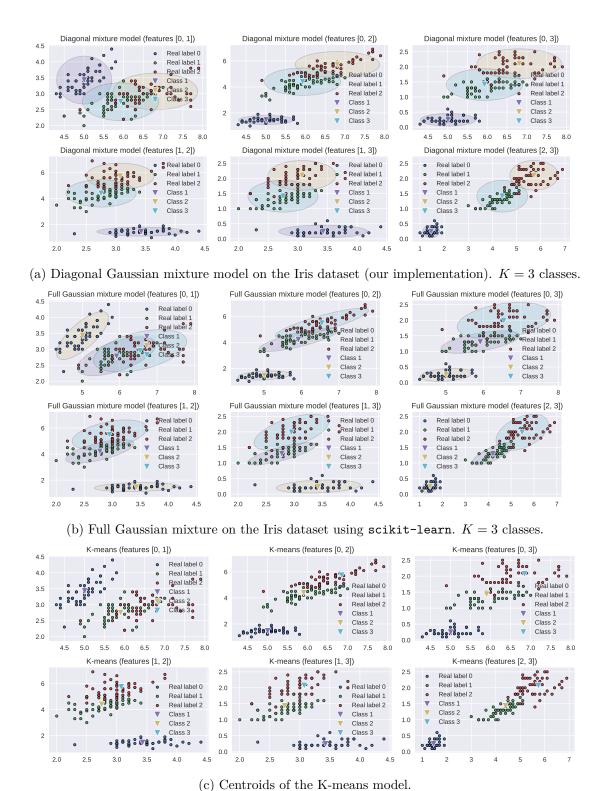


Figure 1: Comparison of the diagonal and full covariance mixture models and K-means for K=3 classes.

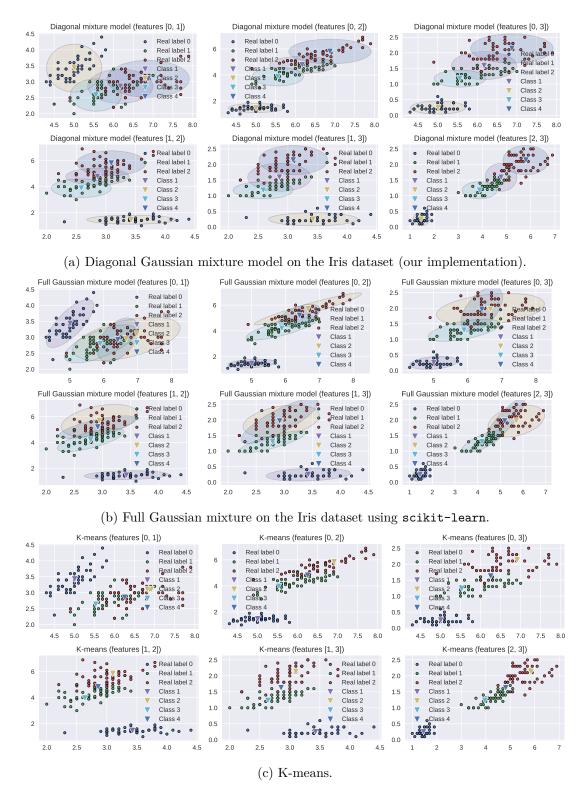
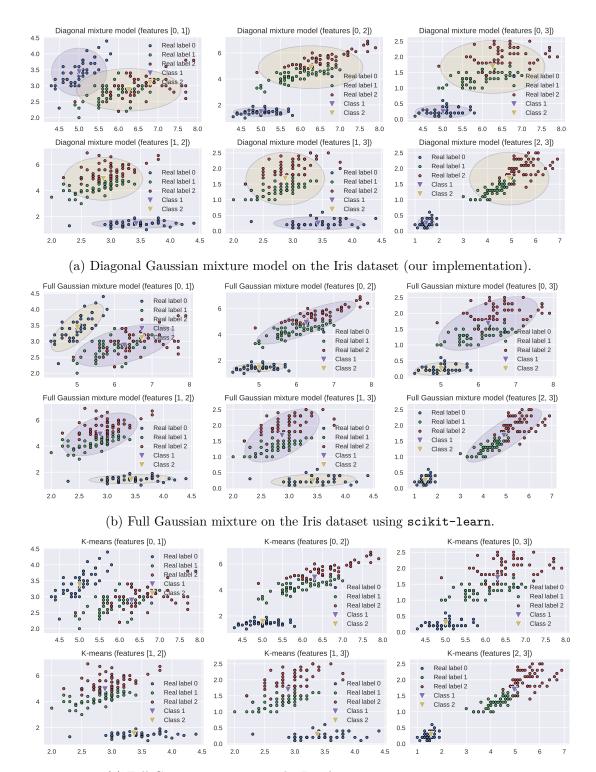
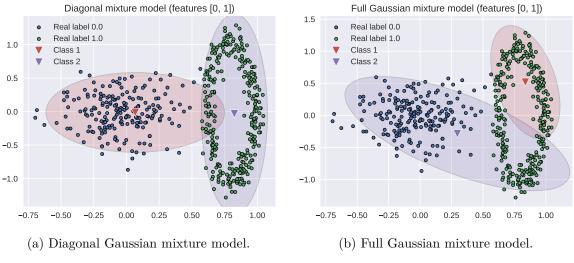


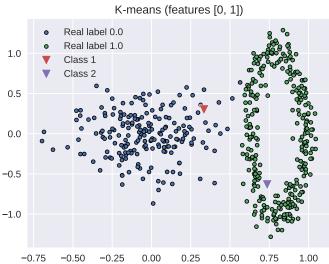
Figure 2: Comparison of the models for K = 4 classes.



(c) Full Gaussian mixture on the Iris dataset using scikit-learn.

Figure 3: Comparison of the models for K=2 classes.





(c) Result of the K-means model: notice how the class centroids are offset from the actuel center of the class.

Figure 4: Comparison of the models on synthetic data where K-means performs worse than mixtures.

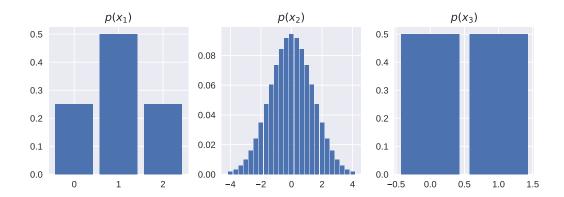


Figure 5: Sanity test: undirected graph with n=3 nodes, and a probability distribution p with factors  $\psi_1=(1,2,1), \ \psi_2(x_2)=\exp(-x_2^2/4), \ \psi_3=(1,1)$  and  $\psi_{i,i+1}=1$  (i.e. independence).

## 2 Graphs, algorithms and Ising

**Question 1.** We recall that for an undirected chain graph G probability distributions factor as

$$p(x) = \frac{1}{Z} \prod_{i=1}^{n} \psi_i(x_i) \prod_{i=1}^{n-1} \psi_{i,i+1}(x_i, x_{i+1})$$
 (5)

The algorithm. The marginal distribution of  $X_i$  can be rewritten as

$$p(x_i) = \frac{1}{Z} \mu_{i-1,i}(x_i) \psi_i(x_i) \mu_{i+1,i}(x_i)$$

where  $\mu_{i-1,i}, \mu_{i+1,i}$  are (forward, backward) messages from i-1 to i and i+1 to i. They are propagated as:

$$\mu_{i,i+1}(x_{i+1}) = \sum_{x} \psi_i(x_i)\psi_{i,i+1}(x_i, x_{i+1})\mu_{i-1,i}(x_i)$$
(6a)

$$\mu_{i,i-1}(x_{i-1}) = \sum_{x_i} \psi_i(x_i) \psi_{i-1,i}(x_{i-1}, x_i) \mu_{i+1,i}(x_i)$$
(6b)

**Practical implementation.** If the state space  $\mathcal{X}$  of the variables  $X_1, \ldots, X_n$  (for instance for binary variables) is discrete we can represent the input functions  $\psi_i$  and  $\psi_{i,i+1}$  as arrays. If not (continuous variables for instance), we can discretize a grid over  $\mathcal{X}$  and precompute an array of values for the factors. Denoting the arrays in bold letters, we forward-propagate by

$$\boldsymbol{\mu}_{i,i+1} = (\boldsymbol{\mu}_{i-1,i} \odot \boldsymbol{\psi}_i) \boldsymbol{\psi}_{i,i+1}$$

and back-propagate by

$$\boldsymbol{\mu}_{i,i-1} = (\boldsymbol{\mu}_{i,i+1} \odot \boldsymbol{\psi}_i) \boldsymbol{\psi}_{i-1,i}^T$$

which allows to compute the marginal distributions using vectorized operations.

Figure 5 shows an implementation for independent edges. The expected marginals are  $X_1 \sim \mathcal{M}(1; 1/4, 1/2, 1/4), X_2 \sim \mathcal{N}(0, 2)$  and  $X_3 \sim \mathcal{B}(1/2)$ .

Question 2. The vertex set V of the graph G=(V,E) induced by the grid has vertices of the form v=(j,k) where  $1 \le j \le w$  and  $1 \le k \le h$ . An easy junction tree to extract from G is given by collapsing the rows into supernodes: the resulting tree T has vertex set  $V_T=\{c_k\}$  where  $c_k=\{(j,k)\}_j$  and looking at the resulting edges shows T is actually an undirected chain  $c_1-c_2-\cdots-c_h$ . Computationally, since there are w=10 columns and h=100 rows, this leads to a reasonably-sized state space (2<sup>10</sup> different states) for each supernode.

$$p(x) = \frac{1}{Z} \prod_{k} \phi_k(x_{c_k}) \prod_{k} \phi_{k,k+1}(x_{c_k}, x_{c_{k+1}})$$

where

$$\phi_k(x_{c_k}) = \prod_{j=1}^w e^{\alpha x_{(j,k)}} \prod_{j=1}^{w-1} e^{\beta \mathbb{1}(x_{(j,k)} = x_{(j+1,k)})}$$
$$\phi_{k,k+1}(x_{c_k}, x_{c_{k+1}}) = \prod_{j=1}^w e^{\beta \mathbb{1}(x_{(j,k)} = x_{(j,k+1)})}$$

The vector  $\phi$  is of size  $2^{10}$ .