

Extending the Bayesian Deep Learning Method MultiSWAG

Scott Brownlie

Master of Science
School of Informatics
University of Edinburgh
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Abstract

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Acknowledgements

Any acknowledgements go here.

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Chapter 1

Introduction

The preliminary material of your report should contain:

- The title page.
- An abstract page.
- Optionally an acknowledgements page.
- The table of contents.

As in this example `skeleton.tex`, the above material should be included between:

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Chapter 2

Background

2.1 Factor Analysis

Factor analysis (FA) is a latent variable model which generates observations $\theta \in \mathbb{R}^D$ as follows. First, a latent vector $\mathbf{h} \in \mathbb{R}^K$, for some $K < D$, is sampled from $p(\mathbf{h}) = \mathcal{N}(\mathbf{0}, \mathbf{I})$. Next, \mathbf{h} is transformed onto a K -dimensional linear subspace of \mathbb{R}^D by left-multiplying it by a *factor loading* matrix $\mathbf{F} \in \mathbb{R}^{D \times K}$. The origin of this subspace is then shifted by adding a bias term $\mathbf{c} \in \mathbb{R}^D$. Finally, the data is perturbed by adding some zero mean Gaussian noise $\varepsilon \in \mathbb{R}^D$ sampled from $\mathcal{N}(\mathbf{0}, \Psi)$, where Ψ is a $D \times D$ diagonal matrix [1]. Putting all this together, an observation $\theta \in \mathbb{R}^D$ is generated according to

$$\theta = \mathbf{F}\mathbf{h} + \mathbf{c} + \varepsilon. \quad (2.1)$$

In the context of this thesis, an observation θ is the parameter vector of a neural network.

It follows that, given \mathbf{h} , the observations θ are Gaussian distributed with mean $\mathbf{F}\mathbf{h} + \mathbf{c}$ and covariance Ψ [1]. Formally,

$$p(\theta|\mathbf{h}) = \mathcal{N}(\mathbf{F}\mathbf{h} + \mathbf{c}, \Psi) = \frac{1}{\sqrt{(2\pi)^D |\Psi|}} \exp\left(-\frac{1}{2}(\theta - \mathbf{F}\mathbf{h} - \mathbf{c})^\top \Psi^{-1}(\theta - \mathbf{F}\mathbf{h} - \mathbf{c})\right), \quad (2.2)$$

where $|\Psi|$ is the *determinant* of Ψ . From [1], integrating $p(\theta|\mathbf{h})$ over \mathbf{h} gives the marginal distribution

$$p(\theta) = \mathcal{N}(\mathbf{c}, \mathbf{F}\mathbf{F}^\top + \Psi). \quad (2.3)$$

The parameters of the model are \mathbf{c} , \mathbf{F} and Ψ . The value of \mathbf{c} which maximises the likelihood of the observed data is the empirical mean of the observations [1], which in this case is θ_{SWA} . Having set the bias term, an expectation-maximisation (EM)

or singular value decomposition (SVD) algorithm can find the maximum likelihood estimates of \mathbf{F} and Ψ [1]. However, both methods require storing all the observations in memory, making them impractical for high-dimensional data, such as the parameter vectors of deep neural networks. Two alternative online algorithms are presented in Chapter 4.

2.2 Bayesian Linear Regression

A linear regression model is a mapping from vector inputs $\mathbf{x} \in \mathbb{R}^D$ to scalar outputs $y \in \mathbb{R}$ via an affine transformation. Given a set of observed input-output pairs, $\mathcal{D} = \{(\mathbf{x}_n, y_n)\}_{n=1}^N$, it is assumed that each output is generated according to

$$y_n = \theta^\top \mathbf{x}_n + \varepsilon \quad (2.4)$$

for some unknown $\theta \in \mathbb{R}^D$. The underlying signal, $\theta^\top \mathbf{x}_n$, is corrupted by additive noise,

$$\varepsilon \sim \mathcal{N}(0, \sigma^2), \quad (2.5)$$

for some $\sigma > 0$ [1]. The model is often written with an explicit bias term, but this can be absorbed into θ by adding a constant of one to the input, leading to the expression in Equation (2.4).

2.2.1 Computing the Posterior

Due to the additive noise, each y_n is a random variable, conditioned on \mathbf{x}_n and θ . Since ε is Gaussian distributed, the conditional pdf of y_n is

$$p(y_n | \theta, \mathbf{x}_n) = \mathcal{N}(\theta^\top \mathbf{x}_n, \sigma^2) = \frac{1}{\sqrt{2\pi\sigma^2}} \exp\left(-\frac{1}{2\sigma^2}(y_n - \theta^\top \mathbf{x}_n)^2\right). \quad (2.6)$$

Assuming that the observations in \mathcal{D} are independent and identically distributed (iid), the log-likelihood of θ having generated the data is

$$\begin{aligned} \log p(\mathcal{D} | \theta) &= \sum_{n=1}^N [\log p(y_n | \theta, \mathbf{x}_n) + \log p(\mathbf{x}_n)] \\ &= \sum_{n=1}^N \left[-\frac{1}{2} \log 2\pi\sigma^2 - \frac{1}{2\sigma^2} (y_n - \theta^\top \mathbf{x}_n)^2 \right] + \sum_{n=1}^N \log p(\mathbf{x}_n) \\ &= -\frac{1}{2\sigma^2} \sum_{n=1}^N (y_n - \theta^\top \mathbf{x}_n)^2 - \frac{N}{2} \log \sigma^2 - \frac{N}{2} \log 2\pi + \sum_{n=1}^N \log p(\mathbf{x}_n) \\ &= -\frac{\beta}{2} \sum_{n=1}^N (y_n - \theta^\top \mathbf{x}_n)^2 + \frac{N}{2} \log \beta + \text{constant}, \end{aligned} \quad (2.7)$$

where $\beta = \frac{1}{\sigma^2}$ [1]. In Bayesian linear regression, a prior distribution for θ is also specified. A common choice is

$$p(\theta) = \mathcal{N}(\mathbf{0}, \alpha^{-1} \mathbf{I}) = \left(\frac{\alpha}{2\pi}\right)^{\frac{D}{2}} \exp\left(-\frac{\alpha}{2} \theta^\top \theta\right) \quad (2.8)$$

for some $\alpha > 0$, which is a hyperparameter known as the *precision* [1]. Applying the logarithm,

$$\begin{aligned} \log p(\theta) &= \frac{D}{2} \log \frac{\alpha}{2\pi} - \frac{\alpha}{2} \theta^\top \theta \\ &= -\frac{\alpha}{2} \theta^\top \theta + \frac{D}{2} \log \alpha + \text{constant}. \end{aligned} \quad (2.9)$$

Now, using Bayes' rule and Equation (2.7) and Equation (2.9), it follows that the log-posterior distribution of θ for fixed α and β is

$$\begin{aligned} \log p(\theta|\mathcal{D}) &= \log \frac{p(\mathcal{D}|\theta)p(\theta)}{p(\mathcal{D})} \\ &= \log p(\mathcal{D}|\theta) + \log p(\theta) - \log p(\mathcal{D}) \\ &= -\frac{\beta}{2} \sum_{n=1}^N (y_n - \theta^\top \mathbf{x}_n)^2 - \frac{\alpha}{2} \theta^\top \theta + \text{constant} \\ &= -\frac{\beta}{2} \sum_{n=1}^N (y_n^2 - 2y_n \theta^\top \mathbf{x}_n + (\theta^\top \mathbf{x}_n)^2) - \frac{\alpha}{2} \theta^\top \theta + \text{constant} \\ &= -\frac{\beta}{2} \sum_{n=1}^N y_n^2 + \beta \sum_{n=1}^N y_n \theta^\top \mathbf{x}_n - \frac{\beta}{2} \sum_{n=1}^N \theta^\top \mathbf{x}_n \mathbf{x}_n^\top \theta - \frac{\alpha}{2} \theta^\top \theta + \text{constant} \\ &= \left(\beta \sum_{n=1}^N y_n \mathbf{x}_n\right)^\top \theta - \frac{1}{2} \theta^\top \left(\alpha \mathbf{I} + \beta \sum_{n=1}^N \mathbf{x}_n \mathbf{x}_n^\top\right) \theta + \text{constant} \\ &= \mathbf{b}^\top \theta - \frac{1}{2} \theta^\top \mathbf{A} \theta + \text{constant}, \end{aligned} \quad (2.10)$$

where

$$\mathbf{A} = \alpha \mathbf{I} + \beta \sum_{n=1}^N \mathbf{x}_n \mathbf{x}_n^\top \quad \text{and} \quad \mathbf{b} = \beta \sum_{n=1}^N y_n \mathbf{x}_n. \quad (2.11)$$

Now, using a result from [1] to complete the square of Equation (2.10),

$$\begin{aligned} \log p(\theta|\mathcal{D}) &= -\frac{1}{2} (\theta - \mathbf{A}^{-1} \mathbf{b})^\top \mathbf{A} (\theta - \mathbf{A}^{-1} \mathbf{b}) + \frac{1}{2} \mathbf{b}^\top \mathbf{A}^{-1} \mathbf{b} + \text{constant} \\ &= -\frac{1}{2} (\theta - \mathbf{A}^{-1} \mathbf{b})^\top \mathbf{A} (\theta - \mathbf{A}^{-1} \mathbf{b}) + \text{constant} \\ &= -\frac{1}{2} (\theta - \mathbf{m})^\top \mathbf{S}^{-1} (\theta - \mathbf{m}) + \text{constant}, \end{aligned} \quad (2.12)$$

where

$$\mathbf{m} = \mathbf{A}^{-1} \mathbf{b} \quad \text{and} \quad \mathbf{S} = \mathbf{A}^{-1}. \quad (2.13)$$

Hence, the posterior distribution of θ is

$$p(\theta|\mathcal{D}) = \mathcal{N}(\mathbf{m}, \mathbf{S}). \quad (2.14)$$

2.3 Gaussian Processes

As in Section 2.2, let $\mathcal{D} = \{(\mathbf{x}_n, y_n)\}_{n=1}^N$ be a set of training examples, where each $\mathbf{x}_n \in \mathbb{R}^D$ and $y_n \in \mathbb{R}$. Considering only the training set for now, a function f can be specified by its values $\{f_n = f(\mathbf{x}_n)\}_{n=1}^N$. Letting $\mathbf{f} = (f_1, \dots, f_N)^\top$, the function f can be seen to generate a sample from a multivariate distribution. If this distribution is constrained to be Gaussian, then it is completely specified by the pdf

$$p(\mathbf{f}) = \mathcal{N}(\mu, \Sigma) \quad (2.15)$$

with parameters $\mu \in \mathbb{R}^D$ and $\Sigma \in \mathbb{R}^{D \times D}$. This is an example of a Gaussian process (GP), which is defined as a collection of random variables such that any finite number of the variables are jointly Gaussian [6]. In this case, the finite subset contains the function values of $\{\mathbf{x}_n\}_{n=1}^N$, but any other set of inputs $\mathbf{x} \in \mathbb{R}^D$ would similarly give rise to a multivariate Gaussian distribution.

The distribution $p(\mathbf{f})$ can be used as a prior on the function values of the training inputs before observing the outputs y_n . In the absence of domain knowledge about which function values are possible, a sensible choice for the prior mean is $\mu = \mathbf{0} \in \mathbb{R}^D$. By definition, the covariance matrix Σ must be positive semi-definite. If Σ was diagonal, function values f_i and f_j , $i, j \in \{1, \dots, N\}$, would be independent, irrespective of how close \mathbf{x}_i is to \mathbf{x}_j . This would lead to extremely noisy functions. To encourage smoothness, it is preferable to define a covariance matrix which assigns greater dependence between function values of nearby inputs. One way to generate covariance matrices with this property is by using the *squared exponential kernel* function [6]. Formally,

$$\Sigma_{i,j} = k(\mathbf{x}_i, \mathbf{x}_j) = \exp\left(-\frac{1}{2}|\mathbf{x}_i - \mathbf{x}_j|^2\right). \quad (2.16)$$

Different *kernel* functions $k(\cdot, \cdot)$ give rise to different covariance matrices, and hence, different priors on the function values.

To simplify the notation in what follows, let $\mathbf{X} \in \mathbb{R}^{N \times D}$ be formed by stacking $\mathbf{x}_1^\top, \dots, \mathbf{x}_N^\top$ in the rows of a matrix. Then the zero mean prior can be written as

$$p(\mathbf{f}) = \mathcal{N}(\mathbf{0}, K(\mathbf{X}, \mathbf{X})), \quad (2.17)$$

where $K(\mathbf{A}, \mathbf{B})$, for $\mathbf{A} \in \mathbb{R}^{N \times D}$ and $\mathbf{B} \in \mathbb{R}^{M \times D}$, is shorthand from [6] for the $N \times M$ matrix with elements

$$K(\mathbf{A}, \mathbf{B})_{i,j} = k(\mathbf{a}_i, \mathbf{b}_j), \quad (2.18)$$

where \mathbf{a}_i and \mathbf{b}_j are the i -th and j -th rows of \mathbf{A} and \mathbf{B} , respectively.

2.3.1 Making Predictions with Gaussian Processes

Suppose that the y_n are noisy observations of the true function values at the training points. That is, $y_n = f_n + \epsilon_n$, where ϵ_n is zero mean Gaussian noise with variance σ_y^2 . Letting $\mathbf{y} = (y_1, \dots, y_N)^\top$, the prior on \mathbf{y} is

$$p(\mathbf{y}) = \mathcal{N}(\mathbf{0}, K(\mathbf{X}, \mathbf{X}) + \sigma_y^2 \mathbf{I}). \quad (2.19)$$

Now suppose that $\mathbf{x}_* \in \mathbb{R}^D$ is a test point outside the training set. The prior on $f_* = f(\mathbf{x}_*)$, the function value at the test point, is $p(f_*) = \mathcal{N}(0, k(\mathbf{x}_*, \mathbf{x}_*))$. Hence, the joint distribution of the observed training outputs \mathbf{y} and the test prediction f_* is

$$\begin{bmatrix} \mathbf{y} \\ f_* \end{bmatrix} \sim \mathcal{N}\left(\mathbf{0}, \begin{bmatrix} K(\mathbf{X}, \mathbf{X}) + \sigma_y^2 \mathbf{I} & K(\mathbf{X}, \mathbf{x}_*) \\ K(\mathbf{x}_*, \mathbf{X}) & k(\mathbf{x}_*, \mathbf{x}_*) \end{bmatrix}\right). \quad (2.20)$$

Using a standard property of joint Gaussian distributions from [1], the conditional distribution of f_* given \mathbf{y} is $p(f_*|\mathbf{y}) = \mathcal{N}(m, s^2)$, where

$$m = K(\mathbf{x}_*, \mathbf{X}) (K(\mathbf{X}, \mathbf{X}) + \sigma_y^2 \mathbf{I})^{-1} \mathbf{y} \quad (2.21)$$

and

$$s^2 = k(\mathbf{x}_*, \mathbf{x}_*) - K(\mathbf{x}_*, \mathbf{X}) (K(\mathbf{X}, \mathbf{X}) + \sigma_y^2 \mathbf{I})^{-1} K(\mathbf{X}, \mathbf{x}_*). \quad (2.22)$$

Writing

$$\mathbf{A} = (K(\mathbf{X}, \mathbf{X}) + \sigma_y^2 \mathbf{I})^{-1} \quad \text{and} \quad \mathbf{k}_* = K(\mathbf{X}, \mathbf{x}_*), \quad (2.23)$$

the conditional distribution can be simplified to

$$p(f_*|\mathbf{y}) = \mathcal{N}(\mathbf{k}_*^\top \mathbf{A} \mathbf{y}, k(\mathbf{x}_*, \mathbf{x}_*) - \mathbf{k}_*^\top \mathbf{A} \mathbf{k}_*). \quad (2.24)$$

2.3.2 Neural Network Kernel Function

Let $f(\mathbf{x}) : \mathbb{R}^D \mapsto \mathbb{R}$ be a neural network with a single hidden layer. Neural networks will be explained in more detail in Section 2.4, but for now it is enough to know that the functional form of the single layer model is

$$f(\mathbf{x}) = b + \sum_{h=1}^H v_h \cdot g(\mathbf{w}_h^\top \mathbf{x}), \quad (2.25)$$

where H is the number of *units* in the hidden layer, the $\mathbf{w}_h \in \mathbb{R}^D$ and $v_h, b \in \mathbb{R}$ are learnable parameters and $g(x) : \mathbb{R} \mapsto \mathbb{R}$ is a non-linear *activation* function [3]. To simplify notation, \mathbf{x} has been augmented with a constant of one such that the bias term of each hidden unit is absorbed into the corresponding \mathbf{w}_h .

Let $\theta \in \mathbb{R}^{Dh+h+1}$ denote all the learnable parameters of the neural network. For any two points $\mathbf{x}, \mathbf{x}' \in \mathbb{R}^D$, the covariance of their predictions is

$$\mathbb{E}_{p(\theta)} \left[(f(\mathbf{x}) - \mathbb{E}_{p(\theta)}[f(\mathbf{x})]) (f(\mathbf{x}') - \mathbb{E}_{p(\theta)}[f(\mathbf{x}')]) \right]. \quad (2.26)$$

Assuming that $b \sim \mathcal{N}(0, \sigma_b^2)$ and $v_h \sim \mathcal{N}(0, \sigma_v^2)$ are independent and that the $\mathbf{w}_h \sim p(\mathbf{w}) = \mathcal{N}(\mathbf{0}, \Sigma_{\mathbf{w}})$ are iid for each hidden unit, Equation (2.26) can be written as

$$\sigma_b^2 + H \sigma_v^2 \mathbb{E}_{p(\mathbf{w})} [g(\mathbf{w}^\top \mathbf{x}) g(\mathbf{w}^\top \mathbf{x}')], \quad (2.27)$$

where \mathbf{w} replaces \mathbf{w}_h due to the iid assumption [7]. Moreover, as $H \rightarrow \infty$, the neural network will converge to a GP [7]. Of course, for this to be a valid GP the activation function g must give rise to a positive definite kernel. One suitable activation is the *error function*,

$$\text{erf}(z) = \frac{2}{\sqrt{\pi}} \int_0^z e^{-t^2} dt. \quad (2.28)$$

Setting $g = \text{erf}$ and evaluating $\mathbb{E}_{p(\mathbf{w})} [\text{erf}(\mathbf{w}^\top \mathbf{x}) \text{erf}(\mathbf{w}^\top \mathbf{x}')]$, it is shown in [7] that the kernel function of the resultant neural network is

$$k_{\text{NN}}(\mathbf{x}, \mathbf{x}') = \frac{2}{\pi} \sin^{-1} \left(\frac{2 \mathbf{x}^\top \Sigma_{\mathbf{w}} \mathbf{x}'}{\sqrt{(1 + 2 \mathbf{x}^\top \Sigma_{\mathbf{w}} \mathbf{x})(1 + 2 \mathbf{x}'^\top \Sigma_{\mathbf{w}} \mathbf{x}')}} \right). \quad (2.29)$$

2.4 Deep Neural Networks

Chapter 3

Related Work and Research Questions

Chapter 4

Online Factor Analysis

4.1 Online Stochastic Gradient Ascent

In situations where learning latent variable models with the classic EM algorithm is slow, [1] suggests optimising the log-likelihood of the model parameters via gradient methods. Since FA is a latent variable model, this approach can be applied here.

In this case, the log-likelihood of the parameters \mathbf{F} and Ψ given the observed data, $\theta_1, \dots, \theta_T$, is

$$L(\mathbf{F}, \Psi) = \frac{1}{T} \sum_{t=1}^T \log p(\theta_t | \mathbf{F}, \Psi). \quad (4.1)$$

The partial derivatives of the log-likelihood with respect to \mathbf{F} and Ψ are therefore

$$\nabla_{\mathbf{F}, \Psi} L(\mathbf{F}, \Psi) = \frac{1}{T} \sum_{t=1}^T \nabla_{\mathbf{F}, \Psi} \log p(\theta_t | \mathbf{F}, \Psi). \quad (4.2)$$

Computing $\nabla_{\mathbf{F}, \Psi} L(\mathbf{F}, \Psi)$ in full would require a pass over all observations, $\theta_1, \dots, \theta_T$. However, a stochastic gradient algorithm can be used instead, as long as the expectation of the sample derivatives is proportional to $\nabla_{\mathbf{F}, \Psi} L(\mathbf{F}, \Psi)$. By using $\nabla_{\mathbf{F}, \Psi} \log p(\theta_t | \mathbf{F}, \Psi)$, $t = 1, \dots, T$, as the sample derivatives, this condition clearly holds. Hence, as long as the partial derivatives $\nabla_{\mathbf{F}, \Psi} \log p(\theta_t | \mathbf{F}, \Psi)$ can be computed efficiently, they can be used in conjunction with SGD or any of its variants to optimise \mathbf{F} and Ψ online.

By adapting an argument for general latent variable models in [1] to FA, the re-

quired sample derivatives can be written as

$$\begin{aligned}
\nabla_{\mathbf{F}, \Psi} \log p(\theta_t | \mathbf{F}, \Psi) &= \frac{1}{p(\theta_t | \mathbf{F}, \Psi)} \nabla_{\mathbf{F}, \Psi} p(\theta_t | \mathbf{F}, \Psi) \\
&= \frac{1}{p(\theta_t | \mathbf{F}, \Psi)} \nabla_{\mathbf{F}, \Psi} \int_{\mathbf{h}_t} p(\theta_t, \mathbf{h}_t | \mathbf{F}, \Psi) \\
&= \frac{1}{p(\theta_t | \mathbf{F}, \Psi)} \int_{\mathbf{h}_t} \nabla_{\mathbf{F}, \Psi} p(\theta_t, \mathbf{h}_t | \mathbf{F}, \Psi) \\
&= \frac{1}{p(\theta_t | \mathbf{F}, \Psi)} \int_{\mathbf{h}_t} p(\theta_t, \mathbf{h}_t | \mathbf{F}, \Psi) \nabla_{\mathbf{F}, \Psi} \log p(\theta_t, \mathbf{h}_t | \mathbf{F}, \Psi) \quad (4.3) \\
&= \int_{\mathbf{h}_t} \frac{p(\theta_t, \mathbf{h}_t | \mathbf{F}, \Psi)}{p(\theta_t | \mathbf{F}, \Psi)} \nabla_{\mathbf{F}, \Psi} \log p(\theta_t, \mathbf{h}_t | \mathbf{F}, \Psi) \\
&= \int_{\mathbf{h}_t} p(\mathbf{h}_t | \theta_t, \mathbf{F}, \Psi) \nabla_{\mathbf{F}, \Psi} \log p(\theta_t, \mathbf{h}_t | \mathbf{F}, \Psi) \\
&= \mathbb{E}_{p(\mathbf{h}_t | \theta_t, \mathbf{F}, \Psi)} [\nabla_{\mathbf{F}, \Psi} \log p(\theta_t, \mathbf{h}_t | \mathbf{F}, \Psi)].
\end{aligned}$$

This is as far as the derivation in [1] goes. However, given the form of the FA model, it is possible to manipulate the sample derivatives further. In particular, using the fact that $\mathbf{h}_t \sim \mathcal{N}(\mathbf{0}, \mathbf{I})$ is independent from \mathbf{F} and Ψ ,

$$\begin{aligned}
\nabla_{\mathbf{F}, \Psi} \log p(\theta_t, \mathbf{h}_t | \mathbf{F}, \Psi) &= \nabla_{\mathbf{F}, \Psi} \log (p(\theta_t | \mathbf{h}_t, \mathbf{F}, \Psi) p(\mathbf{h}_t | \mathbf{F}, \Psi)) \\
&= \nabla_{\mathbf{F}, \Psi} \log (p(\theta_t | \mathbf{h}_t, \mathbf{F}, \Psi) p(\mathbf{h}_t)) \\
&= \nabla_{\mathbf{F}, \Psi} (\log p(\theta_t | \mathbf{h}_t, \mathbf{F}, \Psi) + \log p(\mathbf{h}_t)) \\
&= \nabla_{\mathbf{F}, \Psi} \log p(\theta_t | \mathbf{h}_t, \mathbf{F}, \Psi).
\end{aligned} \quad (4.4)$$

Substituting Equation (4.4) into Equation (4.3),

$$\nabla_{\mathbf{F}, \Psi} \log p(\theta_t | \mathbf{F}, \Psi) = \mathbb{E}_{p(\mathbf{h}_t | \theta_t, \mathbf{F}, \Psi)} [\nabla_{\mathbf{F}, \Psi} \log p(\theta_t | \mathbf{h}_t, \mathbf{F}, \Psi)]. \quad (4.5)$$

Note that $p(\theta_t | \mathbf{h}_t, \mathbf{F}, \Psi)$ is just the Gaussian distribution in Equation (2.2), given \mathbf{F} and Ψ . Hence, substituting $\mathbf{c} = \theta_{\text{SWA}}$ into Equation (2.2) and applying the logarithm,

$$\begin{aligned}
\log p(\theta_t | \mathbf{h}_t, \mathbf{F}, \Psi) &= -\frac{1}{2} (\theta_t - \mathbf{F} \mathbf{h}_t - \theta_{\text{SWA}})^\top \Psi^{-1} (\theta_t - \mathbf{F} \mathbf{h}_t - \theta_{\text{SWA}}) - \frac{1}{2} \log |\Psi| - \frac{D}{2} \log 2\pi \\
&= -\frac{1}{2} (\mathbf{d}_t - \mathbf{F} \mathbf{h}_t)^\top \Psi^{-1} (\mathbf{d}_t - \mathbf{F} \mathbf{h}_t) - \frac{1}{2} \log |\Psi| - \frac{D}{2} \log 2\pi,
\end{aligned} \quad (4.6)$$

where $\mathbf{d}_t = \theta_t - \theta_{\text{SWA}}$. This is convenient, since Equation (4.6) can be differentiated with respect to both \mathbf{F} and Ψ . Of course, this requires access to θ_{SWA} , which is not available during training. As a compromise - and following the approach of the SWAG covariance approximation - θ_{SWA} can be replaced by the running average of the neural network's parameter vectors.

4.1.1 Partial derivatives with respect to \mathbf{F}

From [5], for any symmetric matrix \mathbf{W} ,

$$\nabla_{\mathbf{A}}(\mathbf{x} - \mathbf{A}\mathbf{s})^T \mathbf{W}(\mathbf{x} - \mathbf{A}\mathbf{s}) = -2\mathbf{W}(\mathbf{x} - \mathbf{A}\mathbf{s})\mathbf{s}^T. \quad (4.7)$$

Hence, differentiating Equation (4.6) with respect to \mathbf{F} gives

$$\nabla_{\mathbf{F}} \log p(\theta_t | \mathbf{h}_t, \mathbf{F}, \Psi) = \Psi^{-1}(\mathbf{d}_t - \mathbf{F}\mathbf{h}_t)\mathbf{h}_t^T. \quad (4.8)$$

It then follows from Equation (4.5) that $\nabla_{\mathbf{F}} \log p(\theta_t | \mathbf{F}, \Psi)$ is the expected value of $\Psi^{-1}(\mathbf{d}_t - \mathbf{F}\mathbf{h}_t)\mathbf{h}_t^T$ over the distribution $p(\mathbf{h}_t | \theta_t, \mathbf{F}, \Psi)$. Letting $\mathbb{E}[\cdot]$ denote $\mathbb{E}_{p(\mathbf{h}_t | \theta_t, \mathbf{F}, \Psi)}[\cdot]$ to simplify the notation,

$$\begin{aligned} \nabla_{\mathbf{F}} \log p(\theta_t | \mathbf{F}, \Psi) &= \mathbb{E}[\Psi^{-1}(\mathbf{d}_t - \mathbf{F}\mathbf{h}_t)\mathbf{h}_t^T] \\ &= \Psi^{-1}(\mathbb{E}[\mathbf{d}_t\mathbf{h}_t^T] - \mathbb{E}[\mathbf{F}\mathbf{h}_t\mathbf{h}_t^T]) \\ &= \Psi^{-1}(\mathbf{d}_t\mathbb{E}[\mathbf{h}_t^T] - \mathbf{F}\mathbb{E}[\mathbf{h}_t\mathbf{h}_t^T]). \end{aligned} \quad (4.9)$$

From the E-step of the EM algorithm for FA in [1], $p(\mathbf{h}_t | \theta_t, \mathbf{F}, \Psi) \propto \mathcal{N}(\mathbf{m}_t, \Sigma)$, where

$$\Sigma = (\mathbf{I} + \mathbf{F}^T \Psi^{-1} \mathbf{F})^{-1} \quad \text{and} \quad \mathbf{m}_t = \Sigma \mathbf{F}^T \Psi^{-1} \mathbf{d}_t. \quad (4.10)$$

Hence, using identities from [5],

$$\mathbb{E}[\mathbf{h}_t^T] = \mathbf{m}_t^T \quad \text{and} \quad \mathbb{E}[\mathbf{h}_t\mathbf{h}_t^T] = \Sigma + \mathbf{m}_t\mathbf{m}_t^T. \quad (4.11)$$

Finally, substituting Equation (4.11) into Equation (4.9),

$$\nabla_{\mathbf{F}} \log p(\theta_t | \mathbf{F}, \Psi) = \Psi^{-1}(\mathbf{d}_t\mathbf{m}_t^T - \mathbf{F}(\Sigma + \mathbf{m}_t\mathbf{m}_t^T)). \quad (4.12)$$

4.1.2 Partial derivatives with respect to Ψ

In order to differentiate Equation (4.6) with respect to Ψ , it helps to use the fact that Ψ is diagonal. First consider $\mathbf{X}^{-1} = \text{diag}(\frac{1}{x_1}, \dots, \frac{1}{x_D})$ and $\mathbf{a} = (a_1, \dots, a_D)^T$. Then

$$\mathbf{a}^T \mathbf{X}^{-1} \mathbf{a} = \sum_{d=1}^D \frac{a_d^2}{x_d}, \quad (4.13)$$

and so

$$\frac{\partial}{\partial x_d} \mathbf{a}^T \mathbf{X}^{-1} \mathbf{a} = \frac{-a_d^2}{x_d^2} \quad (4.14)$$

for $d = 1, \dots, D$. Since the partial derivatives of Equation (4.13) with respect to the off-diagonal entries of \mathbf{X} are zero,

$$\begin{aligned}\nabla_{\mathbf{X}}(\mathbf{a}^\top \mathbf{X}^{-1} \mathbf{a}) &= \text{diag}\left(\frac{-a_1^2}{x_1^2}, \dots, \frac{-a_D^2}{x_D^2}\right) \\ &= -\text{diag}\left(\text{diag}(\mathbf{X}^{-2}) \odot \mathbf{a}^2\right),\end{aligned}\quad (4.15)$$

where \odot denotes the element-wise matrix product (with broadcasting, if applicable) and the square of a vector is applied element-wise. Also, when applied to a D -length vector, $\text{diag}(\cdot)$ represents the $D \times D$ diagonal matrix with the vector on its diagonal, and when applied to a $D \times D$ matrix, $\text{diag}(\cdot)$ represents the D -length vector consisting of the diagonal entries of the matrix. Substituting $\mathbf{X} = \Psi$ and $\mathbf{a} = \mathbf{d}_t - \mathbf{F}\mathbf{h}_t$, into Equation (4.15),

$$\nabla_{\Psi}(\mathbf{d}_t - \mathbf{F}\mathbf{h}_t)^\top \Psi^{-1}(\mathbf{d}_t - \mathbf{F}\mathbf{h}_t) = -\text{diag}\left(\text{diag}(\Psi^{-2}) \odot (\mathbf{d}_t - \mathbf{F}\mathbf{h}_t)^2\right). \quad (4.16)$$

Also, using the identity $\nabla_{\mathbf{X}} \log |\mathbf{X}| = \mathbf{X}^{-\top}$ from [5] and the fact that $\Psi^{-\top} = \Psi^{-1}$,

$$\nabla_{\Psi} \log |\Psi| = \Psi^{-1}. \quad (4.17)$$

Hence, the partial derivatives of Equation (4.6) with respect to Ψ are

$$\nabla_{\Psi} \log p(\theta_t | \mathbf{h}_t, \mathbf{F}, \Psi) = \frac{1}{2} \text{diag}\left(\text{diag}(\Psi^{-2}) \odot (\mathbf{d}_t - \mathbf{F}\mathbf{h}_t)^2\right) - \frac{1}{2} \Psi^{-1}. \quad (4.18)$$

Again, letting $\mathbb{E}[\cdot]$ denote $\mathbb{E}_{p(\mathbf{h}_t | \theta_t, \mathbf{F}, \Psi)}[\cdot]$, it follows from Equation (4.5) that

$$\begin{aligned}2 \cdot \nabla_{\Psi} \log p(\theta_t | \mathbf{F}, \Psi) &= \mathbb{E}\left[\text{diag}\left(\text{diag}(\Psi^{-2}) \odot (\mathbf{d}_t - \mathbf{F}\mathbf{h}_t)^2\right) - \Psi^{-1}\right] \\ &= \text{diag}\left(\mathbb{E}\left[\text{diag}(\Psi^{-2}) \odot (\mathbf{d}_t - \mathbf{F}\mathbf{h}_t)^2\right]\right) - \mathbb{E}[\Psi^{-1}] \\ &= \text{diag}\left(\text{diag}(\Psi^{-2}) \odot \mathbb{E}[(\mathbf{d}_t - \mathbf{F}\mathbf{h}_t)^2]\right) - \Psi^{-1}.\end{aligned}\quad (4.19)$$

Now, expanding the expectation inside Equation (4.19) and substituting in the expressions from Equation (4.11),

$$\begin{aligned}\mathbb{E}[(\mathbf{d}_t - \mathbf{F}\mathbf{h}_t)^2] &= \mathbb{E}[\mathbf{d}_t \odot \mathbf{d}_t] - 2\mathbb{E}[\mathbf{d}_t \odot (\mathbf{F}\mathbf{h}_t)] + \mathbb{E}[(\mathbf{F}\mathbf{h}_t) \odot (\mathbf{F}\mathbf{h}_t)] \\ &= \mathbf{d}_t \odot \mathbf{d}_t - 2\mathbf{d}_t \odot (\mathbf{F}\mathbb{E}[\mathbf{h}_t]) + \mathbb{E}[\text{diag}(\mathbf{F}\mathbf{h}_t \mathbf{h}_t^\top \mathbf{F}^\top)] \\ &= \mathbf{d}_t \odot \mathbf{d}_t - 2\mathbf{d}_t \odot (\mathbf{F}\mathbf{m}_t) + \text{diag}(\mathbb{E}[\mathbf{F}\mathbf{h}_t \mathbf{h}_t^\top \mathbf{F}^\top]) \\ &= \mathbf{d}_t \odot \mathbf{d}_t - 2\mathbf{d}_t \odot (\mathbf{F}\mathbf{m}_t) + \text{diag}(\mathbf{F}\mathbb{E}[\mathbf{h}_t \mathbf{h}_t^\top] \mathbf{F}^\top) \\ &= \mathbf{d}_t \odot \mathbf{d}_t - 2\mathbf{d}_t \odot (\mathbf{F}\mathbf{m}_t) + \text{diag}(\mathbf{F}(\Sigma + \mathbf{m}_t \mathbf{m}_t^\top) \mathbf{F}^\top).\end{aligned}\quad (4.20)$$

Finally, substituting Equation (4.20) into Equation (4.19) and rearranging,

$$\begin{aligned} \nabla_{\Psi} \log p(\theta_t | \mathbf{F}, \Psi) = & \text{diag} \left(\frac{1}{2} \text{diag}(\Psi^{-2}) \odot \left(\mathbf{d}_t \odot \mathbf{d}_t - 2\mathbf{d}_t \odot (\mathbf{F}\mathbf{m}_t) \right. \right. \\ & \left. \left. + \text{diag}(\mathbf{F}(\Sigma + \mathbf{m}_t\mathbf{m}_t^T)\mathbf{F}^T) \right) \right) - \frac{1}{2}\Psi^{-1}. \end{aligned} \quad (4.21)$$

4.1.3 Practical implementation

In practice, storing the full $D \times D$ covariance matrix Ψ (or its inverse) would be infeasible for high-dimensional $\theta_t \in \mathbb{R}^D$. However, since Ψ is diagonal and the partial derivatives of the log-likelihood with respect to the off-diagonal are always zero, it suffices to work with the diagonal entries only. All occurrences of $D \times D$ matrices can then be removed from the gradient calculations.

First, note that the partial derivatives with respect to both \mathbf{F} and Ψ depend on \mathbf{m}_t and Σ from Equation (4.10), which themselves depend on Ψ^{-1} . Let $\psi = \text{diag}(\Psi)$ and $\psi^{-n} = \text{diag}(\Psi^{-n})$ for $n \in \mathbb{N}^+$. Then

$$\mathbf{F}^T \Psi^{-1} = (\mathbf{F} \odot \psi^{-1})^T. \quad (4.22)$$

Now, setting $\mathbf{C} = (\mathbf{F} \odot \psi^{-1})^T$ and substituting into Equation (4.10), it follows that

$$\Sigma = (\mathbf{I} + \mathbf{C}\mathbf{F})^{-1} \quad \text{and} \quad \mathbf{m}_t = \Sigma \mathbf{C} \mathbf{d}_t. \quad (4.23)$$

These more efficient values can then be used in Equation (4.12), which itself can be simplified to

$$\nabla_{\mathbf{F}} \log p(\theta_t | \mathbf{F}, \Psi) = \psi^{-1} \odot (\mathbf{d}_t \mathbf{m}_t^T - \mathbf{F}(\Sigma + \mathbf{m}_t \mathbf{m}_t^T)). \quad (4.24)$$

Also, the partial derivatives with respect to Ψ in Equation (4.21) can be re-written with respect to ψ , as

$$\begin{aligned} \nabla_{\Psi} \log p(\theta_t | \mathbf{F}, \Psi) = & \frac{1}{2} \psi^{-2} \odot \left(\mathbf{d}_t \odot \mathbf{d}_t - 2\mathbf{d}_t \odot (\mathbf{F}\mathbf{m}_t) \right. \\ & \left. + \text{sum}((\mathbf{F}(\Sigma + \mathbf{m}_t \mathbf{m}_t^T)) \odot \mathbf{F}, \text{dim} = 1) \right) - \frac{1}{2} \psi^{-1}, \end{aligned} \quad (4.25)$$

where $\text{sum}(\cdot, \text{dim} = 1)$ denotes the operation of summing along the rows of a matrix.

One final point is that the elements of ψ must be positive, since they represent variances. One way to achieve this is by using the re-parameterisation $\psi = \exp \beta$ for

some unconstrained parameter $\beta \in \mathbb{R}^D$. Then the gradient updates can be performed on β instead of ψ . Using the chain rule from calculus,

$$\begin{aligned}\nabla_{\beta} \log p(\theta_t | \mathbf{F}, \Psi) &= \nabla_{\psi} \log p(\theta_t | \mathbf{F}, \Psi) \odot \nabla_{\beta} \exp \beta \\ &= \nabla_{\psi} \log p(\theta_t | \mathbf{F}, \Psi) \odot \exp \beta \\ &= \nabla_{\psi} \log p(\theta_t | \mathbf{F}, \Psi) \odot \psi,\end{aligned}\tag{4.26}$$

where $\nabla_{\psi} \log p(\theta_t | \mathbf{F}, \Psi)$ is given in Equation (4.25).

Pseudo code for online stochastic gradient ascent (SGA) for FA is given in Algorithm 1.

Algorithm 1 Online Stochastic Gradient Ascent for Factor Analysis

Input: Observation dimension D , latent dimension K , learning rate $\alpha > 0$

- 1: Initialise $\mathbf{F} \in \mathbb{R}^{D \times K}$, $\psi = \mathbf{1}^D$, $\bar{\theta}_0 = \mathbf{0}^D$
 - 2: $\beta \leftarrow \log \psi$
 - 3: **for** $t = 1, \dots, T$ **do**
 - 4: Sample observation $\theta_t \in \mathbb{R}^D$
 - 5: $\bar{\theta}_t \leftarrow \bar{\theta}_{t-1} + \frac{1}{t} (\theta_t - \bar{\theta}_{t-1})$
 - 6: $\mathbf{d}_t \leftarrow \theta_t - \bar{\theta}_t$
 - 7: $\mathbf{C} \leftarrow (\mathbf{F} \odot \psi^{-1})^\top$ (with broadcasting)
 - 8: $\Sigma \leftarrow (\mathbf{I} + \mathbf{C}\mathbf{F})^{-1}$
 - 9: $\mathbf{m}_t \leftarrow \Sigma \mathbf{C} \mathbf{d}_t$
 - 10: Compute $\nabla_{\mathbf{F}} \log p(\theta_t | \mathbf{F}, \Psi)$ according to Equation (4.24)
 - 11: Compute $\nabla_{\beta} \log p(\theta_t | \mathbf{F}, \Psi)$ according to Equation (4.26)
 - 12: $\mathbf{F} \leftarrow \mathbf{F} + \alpha \nabla_{\mathbf{F}} \log p(\theta_t | \mathbf{F}, \Psi)$
 - 13: $\beta \leftarrow \beta + \alpha \nabla_{\beta} \log p(\theta_t | \mathbf{F}, \Psi)$
 - 14: $\psi \leftarrow \exp \beta$
 - 15: **end for**
 - 16: $\theta_{\text{SWA}} \leftarrow \bar{\theta}_T$
 - 17: **return** $\theta_{\text{SWA}}, \mathbf{F}, \psi$
-

4.2 Online Expectation-Maximisation

The classic EM algorithm for FA iteratively optimises the log-likelihood of \mathbf{F} and Ψ by alternating “E” and “M” steps until convergence. Using properties of the Kullback-

Leibler divergence, it can be shown that

$$L(\mathbf{F}, \Psi) \geq - \sum_{t=1}^T \mathbb{E}_{q(\mathbf{h}_t|\theta_t)} [\log q(\mathbf{h}_t|\theta_t)] + \sum_{t=1}^T \mathbb{E}_{q(\mathbf{h}_t|\theta_t)} [\log p(\mathbf{h}_t, \theta_t|\mathbf{F}, \Psi)], \quad (4.27)$$

where the first and second terms on the right-hand side are called the *entropy* and *energy*, respectively, and $q(\mathbf{h}_t|\theta_t)$, $t = 1, \dots, T$, are known as the *variational distributions* [1]. The EM algorithm optimises this lower bound on the log-likelihood with respect to \mathbf{F} , Ψ and also $q(\mathbf{h}_t|\theta_t)$, hence the name “variational distributions”. The idea is that, by pushing up the lower bound, the log-likelihood $L(\mathbf{F}, \Psi)$ will hopefully increase as well. In fact, it is guaranteed that each iteration of EM does not decrease $L(\mathbf{F}, \Psi)$, which again follows from the properties of the Kullback-Leibler divergence [1].

4.2.1 E-step

In the batch E-step, \mathbf{F} and Ψ are fixed and Equation (4.27) is maximised with respect to $q(\mathbf{h}_t|\theta_t)$, $t = 1, \dots, T$. From [1], the optimal variational distributions are

$$q(\mathbf{h}_t|\theta_t) = p(\mathbf{h}_t|\theta_t, \mathbf{F}, \Psi) \propto \mathcal{N}(\mathbf{m}_t, \Sigma), \quad (4.28)$$

where \mathbf{m}_t and Σ are given in Equation (4.23). Note that this optimisation can be performed separately for each θ_t as it is sampled, using the estimates of \mathbf{F} and Ψ on iteration t . However, in batch EM all $q(\mathbf{h}_t|\theta_t)$ are updated on every iteration. This is clearly not possible in an online algorithm which discards θ_t before sampling θ_{t+1} . Therefore, as a compromise, in the online version each $q(\mathbf{h}_t|\theta_t)$ will be computed once only, on iteration t , and held fixed thereafter. The only other detail is that the batch algorithm uses θ_{SWA} to compute \mathbf{m}_t . As θ_{SWA} is not available during training, it will be replaced by the running average of the neural network’s parameter vectors, as in the online SGA algorithm from Section 4.1.

4.2.2 M-step

In the batch M-step, $q(\mathbf{h}_t|\theta_t)$, $t = 1, \dots, T$, are fixed and Equation (4.27) is maximised with respect to \mathbf{F} and Ψ . From [1], the optimal values are

$$\mathbf{F} = \mathbf{A}\mathbf{H}^{-1}, \quad (4.29)$$

where

$$\mathbf{A} = \frac{1}{T} \sum_{t=1}^T \mathbf{d}_t \mathbf{m}_t^\top \quad \text{and} \quad \mathbf{H} = \Sigma + \frac{1}{T} \sum_{t=1}^T \mathbf{m}_t \mathbf{m}_t^\top, \quad (4.30)$$

and

$$\Psi = \text{diag} \left(\text{diag} \left(\frac{1}{T} \sum_{t=1}^T \mathbf{d}_t \mathbf{d}_t^\top - 2\mathbf{F}\mathbf{A}^\top + \mathbf{F}\mathbf{H}\mathbf{F}^\top \right) \right). \quad (4.31)$$

Note that this optimisation involves summing over $t = 1, \dots, T$. Moreover, on each iteration all components of the sums in Equation (4.30) and Equation (4.31) are updated. As in the E-step, updating all components is not possible in an online algorithm. Therefore, in the online version these sums will be updated incrementally on each iteration.

4.2.3 Practical implementation

Let $\bar{\mathbf{A}}_t$ and $\bar{\mathbf{H}}_t$ be the estimates of \mathbf{A} and \mathbf{H} from Equation (4.30), respectively, after iteration t of online EM. That is,

$$\bar{\mathbf{A}}_t = \frac{1}{t} \sum_{i=1}^t \mathbf{d}_i \mathbf{m}_i^\top \quad \text{and} \quad \bar{\mathbf{H}}_t = \Sigma + \frac{1}{t} \sum_{i=1}^t \mathbf{m}_i \mathbf{m}_i^\top. \quad (4.32)$$

Then the estimates of \mathbf{F} and Ψ on iteration t become

$$\mathbf{F} = \bar{\mathbf{A}}_t \bar{\mathbf{H}}_t^{-1} \quad (4.33)$$

and

$$\Psi = \text{diag} \left(\text{diag} \left(\frac{1}{t} \sum_{i=1}^t \mathbf{d}_i \mathbf{d}_i^\top - 2\mathbf{F}\bar{\mathbf{A}}_t^\top + \mathbf{F}\bar{\mathbf{H}}_t \mathbf{F}^\top \right) \right). \quad (4.34)$$

As in the online SGA algorithm from Section 4.1, it suffices to work with $\psi = \text{diag}(\Psi)$ instead of Ψ . The diagonal entries of Equation (4.34) can be computed more efficiently as

$$\begin{aligned} \psi &= \frac{1}{t} \sum_{i=1}^t \mathbf{d}_i^2 - 2 \cdot \text{sum}(\mathbf{F} \odot \bar{\mathbf{A}}_t, \text{dim} = 1) + \text{sum}((\mathbf{F}\bar{\mathbf{H}}_t) \odot \mathbf{F}, \text{dim} = 1) \\ &= \frac{1}{t} \sum_{i=1}^t \mathbf{d}_i^2 + \text{sum}((\mathbf{F}\bar{\mathbf{H}}_t) \odot \mathbf{F} - 2\mathbf{F} \odot \bar{\mathbf{A}}_t, \text{dim} = 1). \end{aligned} \quad (4.35)$$

All that remains to complete the online algorithm is to define the incremental update rules for $\bar{\mathbf{A}}_t$ and $\bar{\mathbf{H}}_t$. Both are similar, and for $\bar{\mathbf{A}}_t$ the derivation is

$$\begin{aligned}
\bar{\mathbf{A}}_t &= \frac{1}{t} \sum_{i=1}^t \mathbf{d}_i \mathbf{m}_i^\top \\
&= \frac{1}{t} \left(\sum_{i=1}^{t-1} \mathbf{d}_i \mathbf{m}_i^\top + \mathbf{d}_t \mathbf{m}_t^\top \right) \\
&= \frac{1}{t} \left(\frac{t-1}{t-1} \sum_{i=1}^{t-1} \mathbf{d}_i \mathbf{m}_i^\top + \mathbf{d}_t \mathbf{m}_t^\top \right) \\
&= \frac{1}{t} \left(t \cdot \frac{1}{t-1} \sum_{i=1}^{t-1} \mathbf{d}_i \mathbf{m}_i^\top - \frac{1}{t-1} \sum_{i=1}^{t-1} \mathbf{d}_i \mathbf{m}_i^\top + \mathbf{d}_t \mathbf{m}_t^\top \right) \\
&= \frac{1}{t} \left(t \bar{\mathbf{A}}_{t-1} - \bar{\mathbf{A}}_{t-1} + \mathbf{d}_t \mathbf{m}_t^\top \right) \\
&= \bar{\mathbf{A}}_{t-1} + \frac{1}{t} (\mathbf{d}_t \mathbf{m}_t^\top - \bar{\mathbf{A}}_{t-1}).
\end{aligned} \tag{4.36}$$

Pseudo code for online EM for FA is given in Algorithm 2.

4.3 Factors Initialisation

One detail missing from Algorithms 1 and 2 is how \mathbf{F} is initialised. To encourage diverse factors, it makes sense to initialise \mathbf{F} with orthogonal columns. In practice, this can be achieved by first generating a matrix $\mathbf{A} \in \mathbb{R}^{D \times K}$, whose entries are independently sampled from a standard normal distribution, and then computing its reduced QR decomposition¹. This decomposition is $\mathbf{A} = \mathbf{Q}\mathbf{R}$, where $\mathbf{Q} \in \mathbb{R}^{D \times K}$ is orthogonal and $\mathbf{R} \in \mathbb{R}^{K \times K}$ is upper triangular. Then set $\mathbf{F} = \mathbf{Q}$.

¹<https://pytorch.org/docs/1.9.0/generated/torch.linalg.qr.html>

Algorithm 2 Online Expectation-Maximisation for Factor Analysis

Input: Observation dimension D , latent dimension K

- 1: Initialise $\mathbf{F} \in \mathbb{R}^{D \times K}$, $\psi = \mathbf{1}^D$
 - 2: Initialise $\bar{\theta}_0 = \mathbf{0}^D$, $\bar{\mathbf{A}}_0 = \mathbf{0}^{D \times K}$, $\bar{\mathbf{B}}_0 = \mathbf{0}^{K \times K}$, $\bar{\mathbf{d}}^2_0 = \mathbf{0}^D$
 - 3: **for** $t = 1, \dots, T$ **do**
 - 4: Sample observation $\theta_t \in \mathbb{R}^D$
 - 5: $\bar{\theta}_t \leftarrow \bar{\theta}_{t-1} + \frac{1}{t}(\theta_t - \bar{\theta}_{t-1})$
 - 6: $\mathbf{d}_t \leftarrow \theta_t - \bar{\theta}_t$
 - 7: $\mathbf{C} \leftarrow (\mathbf{F} \odot \psi^{-1})^\top$ (with broadcasting)
 - 8: $\Sigma \leftarrow (\mathbf{I} + \mathbf{C}\mathbf{F})^{-1}$
 - 9: $\mathbf{m}_t \leftarrow \Sigma \mathbf{C} \mathbf{d}_t$
 - 10: $\bar{\mathbf{B}}_t \leftarrow \bar{\mathbf{B}}_{t-1} + \frac{1}{t}(\mathbf{m}_t \mathbf{m}_t^\top - \bar{\mathbf{B}}_{t-1})$
 - 11: $\bar{\mathbf{H}}_t \leftarrow \Sigma + \bar{\mathbf{B}}_t$
 - 12: $\bar{\mathbf{A}}_t \leftarrow \bar{\mathbf{A}}_{t-1} + \frac{1}{t}(\mathbf{d}_t \mathbf{m}_t^\top - \bar{\mathbf{A}}_{t-1})$
 - 13: $\mathbf{F} \leftarrow \bar{\mathbf{A}}_t \bar{\mathbf{H}}_t^{-1}$
 - 14: $\bar{\mathbf{d}}^2_t \leftarrow \bar{\mathbf{d}}^2_{t-1} + \frac{1}{t}(\mathbf{d}_t^2 - \bar{\mathbf{d}}^2_{t-1})$
 - 15: $\psi \leftarrow \bar{\mathbf{d}}^2_t + \text{sum}((\mathbf{F} \bar{\mathbf{H}}_t) \odot \mathbf{F} - 2\mathbf{F} \odot \bar{\mathbf{A}}_t, \text{dim} = 1)$
 - 16: **end for**
 - 17: $\theta_{\text{SWA}} \leftarrow \bar{\theta}_T$
 - 18: **return** $\theta_{\text{SWA}}, \mathbf{F}, \psi$
-

Chapter 5

Online Factor Analysis Experiments

5.1 Methodology

The aim of these experiments is to test how well the online FA algorithms from Section 2.1 are able to fit observations sampled from actual FA models. The form of a FA model is given in Equation (2.3). In particular, it has parameters \mathbf{c} , \mathbf{F} and Ψ . The maximum likelihood estimate of \mathbf{c} , given the data, is just the empirical mean of the sampled observations [1]. This is computed exactly in both online FA algorithms (Algorithms 1 and 2) by maintaining a running average of the observations. The other parameters, \mathbf{F} and Ψ , appear in the FA model as part of the full covariance matrix, $\mathbf{F}\mathbf{F}^\top + \Psi$. Note that right-multiplying \mathbf{F} by any orthogonal matrix $\mathbf{A} \in \mathbb{R}^{K \times K}$ will result in the exact same covariance [1], since

$$(\mathbf{F}\mathbf{A})(\mathbf{F}\mathbf{A})^\top + \Psi = \mathbf{F}\mathbf{A}\mathbf{A}^\top\mathbf{F}^\top + \Psi = \mathbf{F}\mathbf{F}^\top + \Psi. \quad (5.1)$$

Therefore, the \mathbf{F} estimated by an online FA algorithm cannot be directly compared to the true \mathbf{F} . The comparison must be made between the full covariance matrices.

Since the covariance matrices have shape $D \times D$, memory requirements dictate that D , the observation dimension, cannot be too large. Therefore, values of $D = 100$ and $D = 1000$ were used in all experiments in this section. In all cases the latent dimension was set to $K = 10$, meaning that two different observation to latent ratios were tested: $\frac{D}{K} = 10$ and $\frac{D}{K} = 100$. Since computing the maximum likelihood estimate of $\mathbf{c} \in \mathbb{R}^D$ is trivial, in all experiments the entries of \mathbf{c} were simply sampled independently from a standard normal distribution. Each factor loading matrix \mathbf{F} was generated in such a way that its columns spanned the K -dimensional latent space and the conditioning number of the resultant covariance matrix could be controlled. This was achieved by finding

the first K eigenvectors of a symmetric positive semi-definite matrix $\mathbf{M} \in \mathbb{R}^{D \times D}$, and then setting the columns of \mathbf{F} equal to these eigenvectors multiplied by a vector $\mathbf{s} > \mathbf{0} \in \mathbb{R}^D$ (element-wise). By construction, the K eigenvectors are linearly independent and therefore span the latent space, and scaling them by \mathbf{s} affects the conditioning number of $\mathbf{F}\mathbf{F}^\top$. The conditioning number of a matrix is equal to

$$\max_{i,j} \frac{|\lambda_i|}{|\lambda_j|}, \quad (5.2)$$

where the λ_i and λ_j are eigenvalues of the matrix [3]. This ratio can be controlled by specifying the range of \mathbf{s} , or alternatively, the range of \mathbf{s}^2 , which is called the *spectrum*. The larger the ratio of the upper to lower bound of the spectrum, the larger the ratio of the eigenvalues. In each experiment the spectrum was sampled from a uniform distribution with one of the following ranges: $[1, 10]$, $[1, 100]$ or $[1, 1000]$. Finally, the diagonal entries of Ψ were sampled from a uniform distribution with upper bound equal to the maximum value of \mathbf{s}^2 . This is consistent with the FA assumption that an observation, $\boldsymbol{\theta} = \mathbf{F}\mathbf{h} + \mathbf{c} + \boldsymbol{\varepsilon}$, is generated by corrupting the signal $\mathbf{F}\mathbf{h} + \mathbf{c}$ with some random noise $\boldsymbol{\varepsilon} \sim \mathcal{N}(\mathbf{0}, \Psi)$. The full details of how FA models were generated are given in Algorithm 3.

Algorithm 3 Generate a Factor Analysis Model

Input: Observation dimension D , latent dimension K , spectrum range $[a, b]$

- 1: Generate $\mathbf{c} \in \mathbb{R}^D$ by sampling entries independently from $\mathcal{N}(0, 1)$
 - 2: Generate $\mathbf{A} \in \mathbb{R}^{D \times D}$ by sampling entries independently from $\mathcal{N}(0, 1)$
 - 3: $\mathbf{M} \leftarrow \mathbf{A}\mathbf{A}^\top$
 - 4: Compute the K eigenvectors, $\mathbf{v}_1, \dots, \mathbf{v}_K \in \mathbb{R}^D$, corresponding to the K largest eigenvalues of \mathbf{M}
 - 5: Construct the matrix $\mathbf{V}_K \in \mathbb{R}^{D \times K}$ with columns $\mathbf{v}_1, \dots, \mathbf{v}_K$
 - 6: Generate $\mathbf{s}^2 \in \mathbb{R}^D$ by sampling entries independently from $\mathcal{U}(a, b)$
 - 7: $\mathbf{s} \leftarrow \sqrt{\mathbf{s}^2}$ (square root is applied element-wise)
 - 8: $\mathbf{F} \leftarrow \mathbf{V}_K \odot \mathbf{s}$ (with broadcasting)
 - 9: $s_{\max} \leftarrow \max(\mathbf{s}^2)$ (largest element of \mathbf{s}^2)
 - 10: Generate $\boldsymbol{\psi} \in \mathbb{R}^D$ by sampling entries independently from $\mathcal{U}(0, s_{\max})$
 - 11: $\Psi \leftarrow \text{diag}(\boldsymbol{\psi})$
 - 12: **return** $\mathbf{c}, \mathbf{F}, \Psi$
-

Having generated a FA model, observations were sampled according to Equation (2.1). Using the data, the parameters of the model were then estimated by three sep-

arate FA algorithms: online SGA (Algorithm 1), online EM (Algorithm 2) and the scikit-learn `FactorAnalysis` estimator [4]. The scikit-learn implementation is based on the batch SVD algorithm from [1]. The *randomised* version of SVD was used, which for most applications is “sufficiently precise while providing significant speed gains” [4]. Each algorithm was tested on samples of varying size, from 100 up to 100,000 observations. The latent dimension of each approximate model was set to the true latent dimension K . All other hyperparameters of the scikit-learn algorithm were set to their default values. For both online algorithms, \mathbf{F} was initialised as described in Section 4.3, and online SGA ran with a learning rate of 0.001. Additionally, both online algorithms were allowed a *warm-up* period of 100 iterations during which the running averages, $\bar{\boldsymbol{\theta}}_t, \bar{\mathbf{A}}_t, \bar{\mathbf{B}}_t$ and $\bar{\mathbf{d}}^2_t$, were updated, while \mathbf{F} and Ψ remained fixed (TODO: perhaps include in an appendix the derivation that Michael sent, about how this is equivalent to specifying a prior). Each experiment was repeated ten times. In each trial a different random seed was used for generating the true FA model and also initialising the parameters of the online algorithms.

5.2 Results and Discussion

Figure 5.1 shows the relative distance between the true covariance matrix of each FA model and the corresponding estimated covariance matrix of each learning algorithm, as a function of the number of samples used to fit the approximate models. The distance between two matrices is measured by the Frobenius norm (also sometimes called the Euclidean norm) of the difference between the two matrices. The distance is then divided by the Frobenius norm of the true covariance matrix to get the relative distance.

In the first row of the figure are the results corresponding to FA models whose factor loading matrices were generated with a spectrum range of $[1, 10]$. When $D = 100$ (top-left), online EM approximates the true covariance matrix better than online SGA when the number of samples, T , is less than 20,000. For $T \geq 20,000$, the results are very close and the standard error bars overlap. Moreover, for $T \geq 50,000$, both online algorithms are on par with batch SVD. When $D = 1000$ (top-right), online EM beats SGA for $T \geq 200$ and matches batch SVD when $T = 100,000$. This is in contrast to SGA, which appears to get stuck after $T = 20,000$ and does not decrease the distance any further. In the second row of the figure are the results corresponding to a spectrum range of $[1, 100]$. For both $D = 100$ (middle-left) and $D = 1000$ (middle-right), online EM outperforms SGA for $T \geq 200$. Both plots are similar in the sense that online EM

initially decreases the distance quite dramatically, even to a greater extent than batch SVD, then the rate of improvement slows before it again catches up with batch SVD at $T = 100,000$. In the third row of the figure are the results corresponding to a spectrum range of $[1, 1000]$. In this case, online EM is again superior to SGA. Up until $T = 5000$ the results are similar to those in the second row of the figure. However, in this case online EM does not learn as efficiently as T becomes large and in the end it is not able to match batch SVD.

Figure 5.2 shows the 2-Wasserstein distance (TODO: add citation) between the Gaussian distribution defined by each estimated FA model and the Gaussian distribution defined by the true FA model, for the same combinations of observation dimension, latent dimension, spectrum range and sample size. Since all algorithms are able to compute the mean of the true Gaussian distribution exactly, the Wasserstein distance in this case is just a different measure of distance between the true and estimated covariance matrices. Therefore, the shape of the plots in Figure 5.2 is very similar to the shape of the plots in Figure 5.1, albeit the y-axis scales are different.

In summary, these results suggest that online EM is in general better at fitting FA models than online SGA, for all combinations of observation dimension, latent dimension, spectrum range and sample size. The SGA algorithm could possibly be improved by tuning its learning rate, or even using a learning rate scheduler. However, tuning hyperparameters is often a hassle and the fact that this is not required for online EM is another distinct advantage. Compared to batch SVD, the results achieved by online EM are almost identical when $T = 100,000$, except when the spectrum range is equal to $[1, 1000]$. This case would appear to be more challenging for online EM due to the larger conditioning number of the covariance matrix. It appears to get stuck in a local minimum at around $T = 1000$ and the rate of improvement only starts to increase again after $T = 20,000$. Perhaps online EM would be able to catch batch SVD if more than 100,000 samples were used. However, this would be an unreasonably high number of neural network parameter vectors to sample for the use case in this thesis.

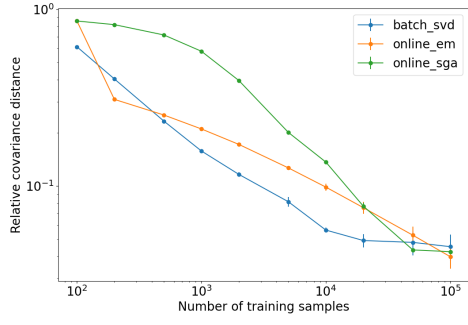
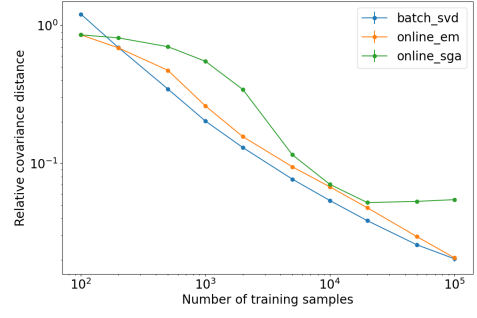
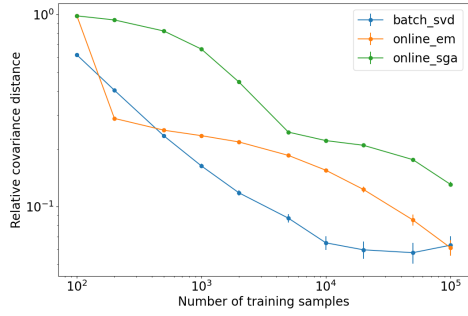
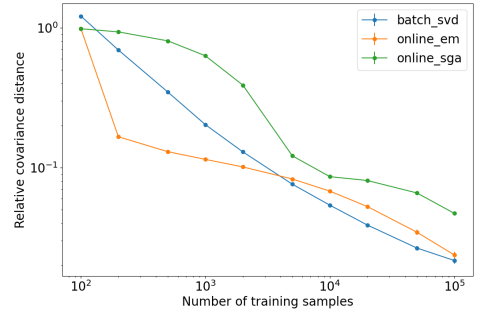
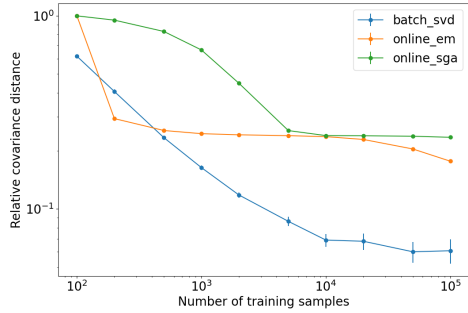
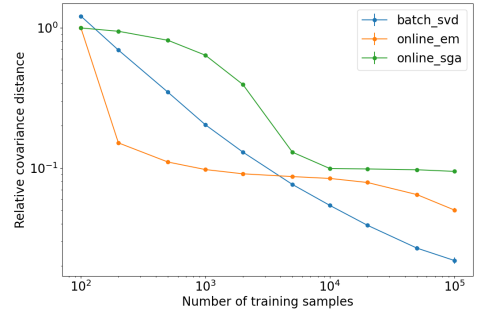
(a) $D = 100, K = 10, \mathbf{s}^2 \sim \mathcal{U}(1, 10)$ (b) $D = 1000, K = 10, \mathbf{s}^2 \sim \mathcal{U}(1, 10)$ (c) $D = 100, K = 10, \mathbf{s}^2 \sim \mathcal{U}(1, 100)$ (d) $D = 1000, K = 10, \mathbf{s}^2 \sim \mathcal{U}(1, 100)$ (e) $D = 100, K = 10, \mathbf{s}^2 \sim \mathcal{U}(1, 1000)$ (f) $D = 1000, K = 10, \mathbf{s}^2 \sim \mathcal{U}(1, 1000)$

Figure 5.1: The relative distance of the estimated FA covariance matrices from the true covariance matrix as a function of the number of samples used to learn the models. The blue, orange and green lines show, for batch SVD, online EM and online SGA, respectively, the Frobenius norm of the difference between the true covariance matrix and the estimated covariance matrix divided by the Frobenius norm of the true covariance matrix. Each data point shows the mean value over ten trials with different random seeds, and standard error bars are also plotted. The different plots correspond to different combinations of the observation dimension D , the latent dimension K and the range of the spectrum \mathbf{s}^2 .

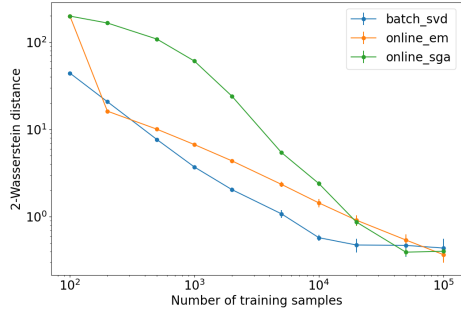
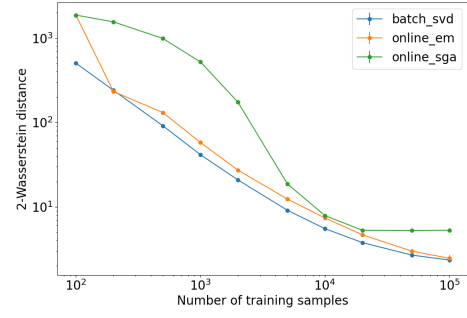
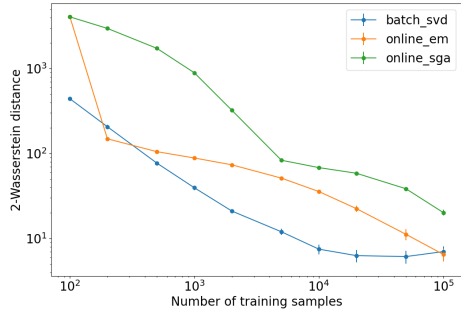
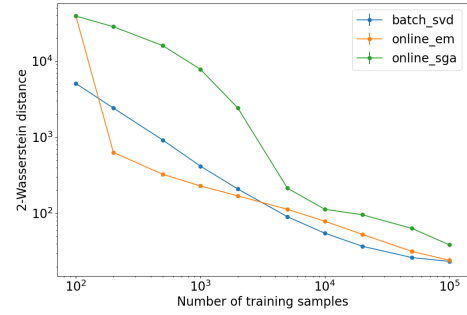
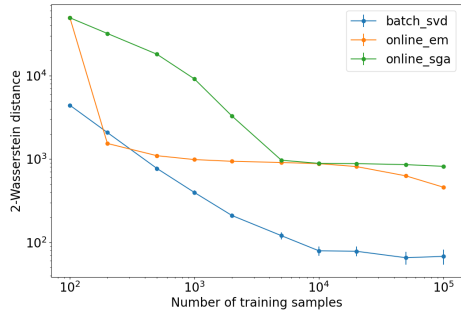
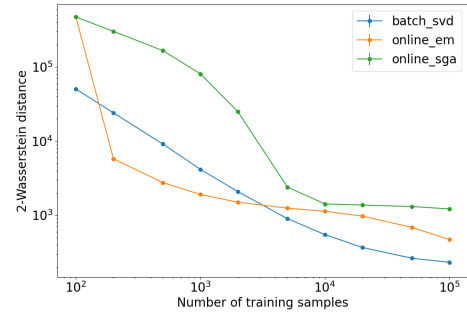
(a) $D = 100, K = 10, \mathbf{s}^2 \sim \mathcal{U}(1, 10)$ (b) $D = 1000, K = 10, \mathbf{s}^2 \sim \mathcal{U}(1, 10)$ (c) $D = 100, K = 10, \mathbf{s}^2 \sim \mathcal{U}(1, 100)$ (d) $D = 1000, K = 10, \mathbf{s}^2 \sim \mathcal{U}(1, 100)$ (e) $D = 100, K = 10, \mathbf{s}^2 \sim \mathcal{U}(1, 1000)$ (f) $D = 1000, K = 10, \mathbf{s}^2 \sim \mathcal{U}(1, 1000)$

Figure 5.2: The 2-Wasserstein distance between the Gaussian distribution defined by each estimated FA model and the Gaussian distribution defined by the true FA model. The blue, orange and green lines show the 2-Wasserstein distance corresponding to batch SVD, online EM and online SGA, respectively. Each data point shows the mean value over ten trials with different random seeds, and standard error bars are also plotted. The different plots correspond to different combinations of the observation dimension D , the latent dimension K and the range of the spectrum \mathbf{s}^2 .

Chapter 6

Linear Regression Experiments

Since the posterior of the linear regression parameter vector can be computed in closed form, the ability of the online FA algorithms to learn the posterior can be evaluated in this case. For a linear regression model trained via SGD, Algorithms 1 and 2 can both be used to fit a Gaussian posterior distribution to the iterates, $\theta_1, \dots, \theta_T$, encountered along the SGD trajectory. The mean and covariance of the learned distributions can then be compared directly to the ground truth. Recall, however, that the posterior distribution in Equation (2.14) refers to specific values of α and β . The parameter β is $\frac{1}{\sigma^2}$, where σ is the standard deviation of the outputs y_n in Equation (2.6). This can be estimated by calculating the empirical standard deviation of $\{y_n\}_{n=1}^N$.

The precision α is a hyperparameter which controls the width of the prior $p(\theta)$. In Equation (2.10), the log-posterior of θ was written as

$$\log p(\theta|\mathcal{D}) = -\frac{\beta}{2} \sum_{n=1}^N (y_n - \theta^\top \mathbf{x}_n)^2 - \frac{\alpha}{2} \theta^\top \theta + \text{constant}. \quad (6.1)$$

Hence, the maximum *a posteriori* (MAP) estimate of θ is that which maximises

$$-\frac{\beta}{2} \sum_{n=1}^N (y_n - \theta^\top \mathbf{x}_n)^2 - \frac{\alpha}{2} \theta^\top \theta. \quad (6.2)$$

Since $\beta > 0$, it is equivalent to minimise this expression scaled by $-\frac{2}{\beta}$, that is,

$$\sum_{n=1}^N (y_n - \theta^\top \mathbf{x}_n)^2 + \frac{\alpha}{\beta} \theta^\top \theta. \quad (6.3)$$

This objective function is analogous to the L2-regularised training loss often used in non-Bayesian linear regression [1], which is

$$\sum_{n=1}^N (y_n - \theta^\top \mathbf{x}_n)^2 + \lambda \theta^\top \theta \quad (6.4)$$

for some $\lambda > 0$. Setting $\lambda = \frac{\alpha}{\beta}$, the optimal θ found by L2-regularised linear regression is the same as the MAP estimate of θ in Bayesian linear regression with prior $p(\theta) = \mathcal{N}(\mathbf{0}, \alpha^{-1}\mathbf{I})$. This means that, given data \mathcal{D} and a particular value of α , the posterior distribution of θ can be estimated by fitting an online FA model to the SGD iterates sampled while minimising Equation (6.4) with $\lambda = \frac{\alpha}{\beta}$.

6.1 Methodology

The aim of these experiments is to test how well the online FA algorithms, fit to the weight vectors sampled during SGD while training a linear regression model, are able to approximate the true posterior distribution of the weights. For this purpose, four real regression datasets from the UCI Machine Learning Repository [2] were used. Namely, the Boston Housing¹, Concrete Compressive Strength², Energy Efficiency³ and Yacht Hydrodynamics⁴ datasets. More details about these datasets are given in Table 6.1. Note that the Energy Efficiency dataset has two target variables, heating load and cooling load, but in these experiments only the heating load was used.

Dataset	No. of Instances	No. of Input Variables
Boston Housing	506	13
Concrete Compressive Strength	1030	8
Energy Efficiency	768	8
Yacht Hydrodynamics	308	6

Table 6.1: UCI regression datasets used in experiments with linear models.

In each experiment, a linear regression model was trained on one of these datasets via mini-batch SGD for 1000 epochs, with a batch size of 32. The L2 regularisation strength was set to $\lambda = \frac{\alpha}{\beta}$. Recall that β is just the reciprocal of the variance of the target variable and α is a hyperparameter which controls the width of the prior on the linear regression weights. Both α and β are used when computing \mathbf{A} in Equation 2.11. In order to balance the effect of the prior term and the data term, α was set to 0.01 times the mean of the diagonal entries of $\beta \sum_{n=1}^N \mathbf{x}_n \mathbf{x}_n^T$.

¹<https://archive.ics.uci.edu/ml/machine-learning-databases/housing>

²<https://archive.ics.uci.edu/ml/datasets/concrete+compressive+strength>

³<https://archive.ics.uci.edu/ml/datasets/energy+efficiency>

⁴<http://archive.ics.uci.edu/ml/datasets/yacht+hydrodynamics>

Before training began, the factors of two FA models were initialised as described in Section 4.3. Their observation dimension was set to the number of weights in the linear regression model, as their purpose was to approximate the posterior distribution of the weights. After an initial warm-up period of training the linear model for 10 epochs, the two FA models began to be updated, one via online SGA and the other via online EM (Algorithms 1 and 2, respectively). An update to each FA model was performed using the weight vector sampled after each mini-batch update to the linear model. As in the experiments in Chapter 5, \mathbf{F} and Ψ remained fixed during the first 100 updates to the FA models, with only the running average being updated in this period. The SGA updates were performed with a learning rate of 0.001.

The mean and covariance of the true posterior of the linear regression weights, for the specific values of α and β , were computed prior to training according to Equation 2.13. After each set of 10 training epochs, the distance between the parameters of the true posterior and the posteriors estimated by the online FA algorithms was computed. For each dataset, this experiment was run with the latent dimension of the FA models set to $K = 1, 2, 3$. Each experiment was repeated ten times. In each trial a different random seed was used for initialising the weights of the linear model and also initialising the parameters of the FA models.

6.2 Results

Chapter 7

Gaussian Processes Experiments

It was shown in Section 2.3 that, under certain conditions on the parameters, a neural network with a single, infinitely wide, hidden layer is equivalent to a GP with kernel function k_{NN} given in Equation (2.29). The kernel depends on $\Sigma_{\mathbf{w}} \in \mathbb{R}^{D \times D}$, which is the covariance of the prior on the weights of each hidden unit. The individual weights of a neural network are usually independent on initialisation, so $\Sigma_{\mathbf{w}}$ can be set to a diagonal matrix $\sigma_w \mathbf{I}$. With this, $k_{\text{NN}}(\mathbf{x}, \mathbf{x}')$ can be evaluated for any two points $\mathbf{x}, \mathbf{x}' \in \mathbb{R}^D$. Hence, given training inputs $\mathbf{X} \in \mathbb{R}^{N \times D}$, training outputs $\mathbf{y} \in \mathbb{R}^N$ and a test point $\mathbf{x}_* \in \mathbb{R}^D$, the Gaussian posterior predictive distribution $p(f_* | \mathbf{X}, \mathbf{y}, \mathbf{x}_*)$ can be computed according to Equation (2.24).

Alternatively, $p(f_* | \mathbf{X}, \mathbf{y}, \mathbf{x}_*)$ can be estimated via an approximate Bayesian model average (BMA). Given a posterior over the parameters of a neural network, $p(\theta | \mathbf{X}, \mathbf{y})$, a parameter vector $\tilde{\theta} \sim p(\theta | \mathbf{X}, \mathbf{y})$ can be used to construct a specific instance of the neural network and make a prediction for \mathbf{x}_* . If this procedure is repeated multiple times, the mean and variance of $p(f_* | \mathbf{X}, \mathbf{y}, \mathbf{x}_*)$ can be estimated and compared directly to the ground truth. To obtain the parameter posterior, the online FA algorithms can be used to fit a Gaussian posterior to the iterates of the parameter vector, $\theta_1, \dots, \theta_T$, encountered while training the single layer neural network via SGD. Of course, in practice it is impossible to train a neural network with an infinitely wide hidden layer, but for these experiments it will suffice to consider neural networks with a few hundred hidden units.

7.1 Methodology

7.2 Results

Chapter 8

Conclusions

8.1 Final Reminder

The body of your dissertation, before the references and any appendices, *must* finish by page 40. The introduction, after preliminary material, should have started on page 1.

You may not change the dissertation format (e.g., reduce the font size, change the margins, or reduce the line spacing from the default 1.5 spacing). Over length or incorrectly-formatted dissertations will not be accepted and you would have to modify your dissertation and resubmit. You cannot assume we will check your submission before the final deadline and if it requires resubmission after the deadline to conform to the page and style requirements you will be subject to the usual late penalties based on your final submission time.

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