## Bayesian learning II

## Bayesian learning II

Bayes decision theory tells us that in this context we should consider the quantity  $\Pr(\omega_i | s, x)$  where the involvement of the training sequence has been made explicit.

$$\begin{split} \Pr(\boldsymbol{\omega}_{i}|\mathbf{s},\mathbf{x}) \;&=\; \sum_{h\in\mathcal{H}} \Pr(\boldsymbol{\omega}_{i},h|\mathbf{s},\mathbf{x}) \\ &=\; \sum_{h\in\mathcal{H}} \Pr(\boldsymbol{\omega}_{i}|h,\mathbf{s},\mathbf{x}) \Pr(h|\mathbf{s},\mathbf{x}) \\ &=\; \sum_{h\in\mathcal{H}} \Pr(\boldsymbol{\omega}_{i}|h,\mathbf{x}) \Pr(h|\mathbf{s}). \end{split}$$

Here we have re-introduced  $\mathcal{H}$  using marginalisation. In moving from line 2 to line 3 we are assuming some independence properties.

## So our classification should be

$$\omega = \operatorname*{argmax}_{\omega \in \{\omega_1,...,\omega_c\}} \sum_{\mathbf{h} \in \mathcal{H}} \Pr(\omega | \mathbf{h}, \mathbf{x}) \Pr(\mathbf{h} | \mathbf{s})$$

If  ${\mathcal H}$  is infinite the sum becomes an integral. So for example for a neural network

 $\boldsymbol{\omega} = \underset{\boldsymbol{\omega} \in \{\boldsymbol{\omega}_1, ..., \boldsymbol{\omega}_c\}}{\operatorname{argmax}} \int_{\mathbb{R}^W} \Pr(\boldsymbol{\omega} | \mathbf{w}, \mathbf{x}) \Pr(\mathbf{w} | \mathbf{s}) \, d\mathbf{w}$ 

where W is the number of weights in w.

## Bayesian learning II

Why might this make any difference? (Aside from the fact that we now know it's optimal!)

Example 1: Say  $|\mathcal{H}| = 3$  and  $h(\mathbf{x}) = Pr(\mathbf{x} \text{ is in class } C_1)$  for a 2 class problem.

$$\begin{aligned} & \Pr(h_1|\mathbf{s}) = 0.4 \\ & \Pr(h_2|\mathbf{s}) = \Pr(h_3|\mathbf{s}) = 0.3 \end{aligned}$$

Now, say we have an x for which

$$\begin{aligned} h_1(\mathbf{x}) &= 1\\ h_2(\mathbf{x}) &= h_3(\mathbf{x}) = 0 \end{aligned}$$

so  $h_{MAP}$  says that x is in class  $C_1$ .

#### Bayesian learning II

However,

$$\begin{aligned} \Pr(\text{class 1}|\mathbf{s}, \mathbf{x}) &= 1 \times 0.4 + 0 \times 0.3 + 0 \times 0.3 \\ &= 0.4 \\ \Pr(\text{class 2}|\mathbf{s}, \mathbf{x}) &= 0 \times 0.4 + 1 \times 0.3 + 1 \times 0.3 \\ &= 0.6 \end{aligned}$$

so class  $C_2$  is the more probable!

In this case the Bayes optimal approach in fact leads to a different answer.

### A more in-depth example

## A more in-depth example

Let's take this a step further and work through something a little more complex in detail. For a two-class classification problem, with h(x) denoting  $Pr(C_1|h, x)$  and  $x \in \mathbb{R}$ :

Hypotheses: We have three hypotheses

$$\begin{split} h_1(x) &= \exp(-(x-1)^2) \\ h_2(x) &= \exp(-(2x-2)^2) \\ h_3(x) &= \exp(-(1/10)(x-3)^2) \end{split}$$

<u>Prior</u>: The prior is  $Pr(h_1) = 0.1$ ,  $Pr(h_2) = 0.05$  and  $Pr(h_3) = 0.85$ .

We see the examples  $(0.5, C_1)$ ,  $(0.9, C_1)$ ,  $(3.1, C_2)$  and  $(3.4, C_1)$ .

 $\underline{Likelihood}$ : For the individual hypotheses the likelihoods are given by

 $Pr(s|h) = h(x_1)h(x_2)[1 - h(x_3)]h(x_4)$ 

Which in this case tells us

 $Pr(s|h_1) = 0.0024001365$  $Pr(s|h_2) = 0.0031069836$  $Pr(s|h_3) = 0.0003387476$ 

Posterior: Multiplying by the priors and normalising gives

 $\begin{aligned} & \Pr(h_1 | \mathbf{s}) = 0.3512575000 \\ & \Pr(h_2 | \mathbf{s}) = 0.2273519164 \\ & \Pr(h_3 | \mathbf{s}) = 0.4213905836 \end{aligned}$ 

A more in-depth example

Now let's classify the point x' = 2.5.

#### We need

 $\begin{aligned} \Pr(C_1|\mathbf{s},\mathbf{x}') &= \Pr(C_1|h_1)\Pr(h_1|\mathbf{s}) + \Pr(C_1|h_2)\Pr(h_2|\mathbf{s}) + \Pr(C_1|h_3)\Pr(h_3|\mathbf{s}) \\ &= 0.6250705317 \end{aligned}$ 

So: it's most likely to be in class  $C_1$ , but not with great certainty.

The Bayesian approach to neural networks

Let's now see how this can be applied to *neural networks*. We have:

- A neural network computing a function  $f(\mathbf{w}; \mathbf{x})$ .
- A training sequence  $s = ((x_1, y_1), \dots, (x_m, y_m))$ , split into

$$\mathbf{y} = (\mathbf{y}_1 \ \mathbf{y}_2 \ \cdots \ \mathbf{y}_m)$$

 $\operatorname{and}$ 

$$\mathbf{X} = (\mathbf{x}_1 \ \mathbf{x}_2 \ \cdots \ \mathbf{x}_m)$$

The prior distribution  $p(\mathbf{w})$  is now on the weight vectors, and Bayes' theorem tells us that

$$p(\mathbf{w}|\mathbf{s}) = p(\mathbf{w}|\mathbf{X}, \mathbf{y}) = \frac{p(\mathbf{y}|\mathbf{w}, \mathbf{X})p(\mathbf{w}|\mathbf{X})}{p(\mathbf{y}|\mathbf{X})}$$

Nothing new so far...

### The Bayesian approach to neural networks

As usual, we don't consider uncertainty in x and so X will be omitted. Consequently

$$p(\mathbf{w}|\mathbf{y}) = \frac{p(\mathbf{y}|\mathbf{w})p(\mathbf{w})}{p(\mathbf{y})}$$

where

$$p(\mathbf{y}) = \int_{\mathbb{R}^W} p(\mathbf{y}|\mathbf{w}) p(\mathbf{w}) d\mathbf{w}$$

 $p(\mathbf{y}|\mathbf{w})$  is a model of the noise corrupting the labels and as previously is the likelihood function.

# Reminder: the general Gaussian density

Reminder: we're going to be making a lot of use of the general Gaussian density  $\mathcal{N}(\mu, \Sigma)$  in d dimensions

$$p(\mathbf{z}) = (2\pi)^{-d/2} |\boldsymbol{\Sigma}|^{-1/2} \exp\left[-\frac{1}{2}\left((\mathbf{z} - \boldsymbol{\mu})^{\mathsf{T}} \boldsymbol{\Sigma}^{-1} (\mathbf{z} - \boldsymbol{\mu})\right)\right]$$

an density,  $\mu = [0 \ 0]$ ,  $\Sigma = I$ 

where  $\mu$  is the mean vector and  $\Sigma$  is the covariance matrix.



## The Bayesian approach to neural networks

 $p(\mathbf{w})$  is typically a *broad distribution* to reflect the fact that in the absence of any data we have little idea of what  $\mathbf{w}$  might be.

When we see some data the above equation tells us how to obtain  $p(\mathbf{w}|\mathbf{y})$ . This will typically be *more localised*.



To put this into practice we need expressions for  $p(\mathbf{w})$  and  $p(\mathbf{y}|\mathbf{w})$ .

## The Gaussian prior

A common choice for  $p(\mathbf{w})$  is the Gaussian prior with zero mean and  $\Sigma = \sigma^2 \mathbf{I}$ 

SO

$$p(\mathbf{w}) = (2\pi)^{-W/2} \sigma^{-W} \exp\left[-\frac{\mathbf{w}^{T} \mathbf{w}}{2\sigma^{2}}\right]$$

Note that  $\sigma$  controls the distribution of other parameters.

- Such parameters are called *hyperparameters*.
- Assume for now that they are both fixed and known.

Hyperparameters can be learnt using s through the application of more advanced techniques.

### The Bayesian approach to neural networks

Physicists like to express quantities such as p(w) in terms of a measure of "energy". The expression is therefore usually re-written as

$$p(\mathbf{w}) = \frac{1}{Z_W(\alpha)} \exp\left(-\frac{\alpha}{2} ||\mathbf{w}||^2\right)$$

where

$$E_W(\mathbf{w}) = \frac{1}{2} ||\mathbf{w}||^2$$
$$Z_W(\alpha) = \left(\frac{2\pi}{\alpha}\right)^{d/2}$$
$$\alpha = \frac{1}{\sigma^2}$$

This is simply a re-arranged version of the more usual equation.

# The Bayesian approach to neural networks

This expression can also be rewritten in physicist-friendly form

$$p(\mathbf{y}|\mathbf{w}) = \frac{1}{Z_{\mathbf{y}}(\beta)} \exp\left(-\beta E_{\mathbf{y}}(\mathbf{w})\right)$$

where

$$\beta = \frac{1}{\sigma_n^2}$$

$$Z_{\mathbf{y}}(\beta) = \left(\frac{2\pi}{\beta}\right)^{m/2}$$

$$E_{\mathbf{y}}(\mathbf{w}) = \frac{1}{2}\sum_{i=1}^m (y_i - f(\mathbf{w}; \mathbf{x}_i))^2$$

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Here,  $\beta$  is a second *hyperparameter*. Again, we assume it is fixed and known, although it can be learnt using s using more advanced techniques.

## The Gaussian noise model for regression

We've already seen that for a regression problem with zero mean Gaussian noise having variance  $\sigma_n^2$ 

$$y_{i} = f(\mathbf{x}_{i}) + \epsilon_{i}$$
$$p(\epsilon_{i}) = \frac{1}{\sqrt{2\pi\sigma_{n}^{2}}} \exp\left(-\frac{\epsilon_{i}^{2}}{2\sigma_{n}^{2}}\right)$$

where f corresponds to some unknown network, the likelihood function is

$$p(\mathbf{y}|\mathbf{w}) = \frac{1}{(2\pi\sigma_n^2)^{m/2}} \exp\left(-\frac{1}{2\sigma_n^2} \sum_{i=1}^m (y_i - f(\mathbf{w}; \mathbf{x}_i))^2\right)$$

Note that there are now two variances:  $\sigma^2$  for the prior and  $\sigma_n^2$  for the noise.

# The Bayesian approach to neural networks

Combining the two boxed equations gives

$$p(\mathbf{w}|\mathbf{y}) = \frac{1}{Z_S(\alpha, \beta)} \exp(-S(\mathbf{w}))$$

where

$$S(\mathbf{w}) = \alpha E_W(\mathbf{w}) + \beta E_y(\mathbf{w})$$

The quantity

$$Z_{S}(\alpha,\beta) = \int_{\mathbb{R}^{W}} \exp(-S(\mathbf{w})) d\mathbf{w}$$

normalises the density. Recall that this is called the *evidence*.

## Example I: gradient descent revisited...

To find  $h_{\text{MAP}}$  (in this scenario by finding  $\mathbf{w}_{\text{MAP}})$  we therefore maximise

$$p(\mathbf{w}|\mathbf{y}) = \frac{1}{Z_{S}(\alpha, \beta)} \exp(-(\alpha E_{W}(\mathbf{w}) + \beta E_{y}(\mathbf{w})))$$

or equivalently find

$$\mathbf{w}_{\text{MAP}} = \underset{\mathbf{w}}{\text{argmin}} \frac{\alpha}{2} ||\mathbf{w}||^2 + \frac{\beta}{2} \sum_{i=1}^{m} (y_i - f(\mathbf{w}; \mathbf{x}_i))^2$$

This algorithm has also been used a lot in the neural network literature and is called the *weight decay* technique.



Example II: two-class classification in two dimensions

# The Bayesian approach to neural networks

What happens as the number m of examples increases?

- The first term *corresponding to the prior* remains fixed.
- The second term *corresponding to the likelihood* increases.

So for small training sequences the prior dominates, but for large ones  $h_{\rm ML}$  is a good approximation to  $h_{\rm MAP}.$ 

### The Bayesian approach to neural networks

Where have we got to ...? We have obtained

$$p(\mathbf{w}|\mathbf{y}) = \frac{1}{Z_{S}(\alpha, \beta)} \exp(-(\alpha E_{W}(\mathbf{w}) + \beta E_{y}(\mathbf{w}))$$
$$Z_{S}(\alpha, \beta) = \int_{\mathbb{R}^{W}} \exp(-(\alpha E_{W}(\mathbf{w}) + \beta E_{y}(\mathbf{w}))) d\mathbf{w}$$

Translating the expression for the *Bayes optimal* solution given on the first slide of this handout into the current scenario, we need to compute

$$p(Y|\mathbf{y}, \mathbf{x}) = \int_{\mathbb{R}^W} p(\mathbf{y}|\mathbf{w}, \mathbf{x}) p(\mathbf{w}|\mathbf{y}) \, d\mathbf{w}$$

Easy huh? Unfortunately not...

### The Bayesian approach to neural networks

In order to make further progress it's therefore necessary to perform integrals of the general form

$$_{W}F(\mathbf{w})p(\mathbf{w}|\mathbf{y})d\mathbf{w}$$

for various functions  ${\sf F}$  and this is generally not possible.

There are two ways to get around this:

1. We can use an *approximate form* for  $p(\mathbf{w}|\mathbf{y})$ .

2. We can use *Monte Carlo* methods.

# Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$

The first approach introduces a Gaussian approximation to  $p(\mathbf{w}|\mathbf{y})$  by using a Taylor expansion of

$$S(\mathbf{w}) = \alpha E_W(\mathbf{w}) + \beta E_y(\mathbf{w})$$

at  $\mathbf{w}_{MAP}$ .

This allows us to use a standard integral.

The result will be *approximate* but we hope it's good!

Let's recall how Taylor series work...

## Reminder: Taylor expansion

In one dimension the Taylor expansion about a point  $x_0\in\mathbb{R}$  for a function  $f:\mathbb{R}\to\mathbb{R}$  is

$$f(x) \approx f(x_0) + \frac{1}{1!}(x - x_0)f'(x_0) + \frac{1}{2!}(x - x_0)^2 f''(x_0) + \dots + \frac{1}{k!}(x - x_0)^k f^k(x_0)$$

What does this look like for the kinds of function we're interested in? We can try to approximate

$$\exp\left(-f(x)\right)$$

where

$$f(\mathbf{x}) = \mathbf{x}^4 - \frac{1}{2}\mathbf{x}^3 - 7\mathbf{x}^2 - \frac{5}{2}\mathbf{x} + 22$$

This has a form similar to S(w), but in one dimension.

### Reminder: Taylor expansion

The functions of interest look like this:



By replacing -f(x) with its Taylor expansion about its maximum, which is at

$$x_{max} = 2.1437$$

we can see what the approximation to  $\exp(-f(x))$  looks like. Note that the exp hugely emphasizes peaks.

#### Reminder: Taylor expansion

Here are the approximations for k = 1, k = 2 and k = 3.





# Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$

Applying this to  $S(\mathbf{w})$  and expanding around  $\mathbf{w}_{MAP}$ 

$$\begin{split} S(\mathbf{w}) &\approx S(\mathbf{w}_{\text{MAP}}) + (\mathbf{w} - \mathbf{w}_{\text{MAP}})^{\mathsf{T}} \nabla S(\mathbf{w})|_{\mathbf{w}_{\text{MAP}}} \\ &+ \frac{1}{2} (\mathbf{w} - \mathbf{w}_{\text{MAP}})^{\mathsf{T}} \mathbf{A} (\mathbf{w} - \mathbf{w}_{\text{MAP}}) \end{split}$$

notice the following:

- As  $\mathbf{w}_{MAP}$  minimizes the function the first derivatives are zero and the corresponding term in the Taylor expansion disappears.
- The quantity  $\mathbf{A} = \nabla \nabla S(\mathbf{w})|_{\mathbf{w}_{MAP}}$  can be simplified.

This is because

$$\begin{split} \mathbf{A} &= \nabla \nabla (\alpha \mathsf{E}_W(\mathbf{w}) + \beta \mathsf{E}_\mathbf{y}(\mathbf{w})) |_{\mathbf{w}_{MAP}} \\ &= \alpha \mathbf{I} + \beta \nabla \nabla \mathsf{E}_\mathbf{y}(\mathbf{w}_{MAP}) \end{split}$$

### Reminder: Taylor expansion

In multiple dimensions the Taylor expansion for k = 2 is

$$f(\mathbf{x}) \approx f(\mathbf{x}_0) + \frac{1}{1!} (\mathbf{x} - \mathbf{x}_0)^T \nabla f(\mathbf{x}) \big|_{\mathbf{x}_0} + \frac{1}{2!} (\mathbf{x} - \mathbf{x}_0)^T \nabla^2 f(\mathbf{x}_0) \big|_{\mathbf{x}_0} (\mathbf{x} - \mathbf{x}_0)$$

where  $\nabla$  denotes gradient

$$abla f(\mathbf{x}) = \left( \begin{array}{ccc} \frac{\partial f(\mathbf{x})}{\partial x_1} & \frac{\partial f(\mathbf{x})}{\partial x_2} & \cdots & \frac{\partial f(\mathbf{x})}{\partial x_n} \end{array} \right)$$

and  $\nabla^2 f(\mathbf{x})$  is the matrix with elements

 $M_{ij} = \frac{\partial^2 f(\mathbf{x})}{\partial x_i \partial x_j}$ 

(Although this looks complicated, it's just the obvious extension of the 1-dimensional case.)

### Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$

Defining

$$\Delta \mathbf{w} = \mathbf{w} - \mathbf{w}_{MAP}$$

we now have

$$S(\mathbf{w}) \approx S(\mathbf{w}_{MAP}) + \frac{1}{2} \Delta \mathbf{w}^{T} \mathbf{A} \Delta \mathbf{w}$$

The vector  $\mathbf{w}_{MAP}$  can be obtained using any standard optimization method (such as *backpropagation*).

The quantity  $\nabla \nabla E_{\mathbf{y}}(\mathbf{w})$  can be evaluated using an extended form of backpropagation.

### A useful integral

Dropping for this slide only the special meanings usually given to vectors  $\mathbf{x}$  and  $\mathbf{y}$ , here is a useful standard integral:

If  $\mathbf{A} \in \mathbb{R}^{n \times n}$  is symmetric then for  $\mathbf{b} \in \mathbb{R}^n$  and  $\mathbf{c} \in \mathbb{R}$ 

$$\begin{split} \int_{\mathbb{R}^n} \exp\left(-\frac{1}{2}\left(\mathbf{x}^\mathsf{T}\mathbf{A}\mathbf{x} + \mathbf{x}^\mathsf{T}\mathbf{b} + \mathbf{c}\right)\right) d\mathbf{x} \\ &= (2\pi)^{n/2} |\mathbf{A}|^{-1/2} \exp\left(-\frac{1}{2}\left(\mathbf{c} - \frac{\mathbf{b}^\mathsf{T}\mathbf{A}^{-1}\mathbf{b}}{4}\right)\right) \end{split}$$

At the beginning of the course, two exercises were set involving the evaluation of this integral.

To make this easy to refer to, let's call it the BIG INTEGRAL.

# Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$

We now have

$$p(\mathbf{w}|\mathbf{y}) \approx \frac{1}{Z(\alpha, \beta)} \exp\left(-S(\mathbf{w}_{MAP}) - \frac{1}{2}\Delta \mathbf{w}^{\mathsf{T}} \mathbf{A} \Delta \mathbf{w}\right)$$

where  $\Delta \mathbf{w} = \mathbf{w} - \mathbf{w}_{\text{MAP}}$  and using the *BIG INTEGRAL* 

$$Z(\alpha,\beta) = (2\pi)^{W/2} |\mathbf{A}|^{-1/2} \exp(-S(\mathbf{w}_{MAP}))$$

Our earlier discussion tells us that given a new input  $\mathbf{x}$  we should calculate

$$p(\mathbf{Y}|\mathbf{y},\mathbf{x}) = \int_{\mathbb{R}^W} p(\mathbf{y}|\mathbf{w},\mathbf{x}) p(\mathbf{w}|\mathbf{y}) d\mathbf{w}$$

p(y|w, x) is just the *likelihood* so...

# Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$

The likelihood we're using is

$$p(\mathbf{y}|\mathbf{w}, \mathbf{x}) = \frac{1}{\sqrt{2\pi\sigma^2}} \exp\left(-\frac{(\mathbf{y} - f(\mathbf{w}; \mathbf{x}))^2}{2\sigma^2}\right)$$
$$\propto \exp\left(-\frac{\beta}{2}(\mathbf{y} - f(\mathbf{w}; \mathbf{x}))^2\right)$$

and plugging it into the integral gives

$$p(\mathbf{y}|\mathbf{x},\mathbf{y}) \propto \int_{\mathbb{R}^W} \exp\left(-\frac{\beta}{2}(\mathbf{y} - f(\mathbf{w};\mathbf{x}))^2\right) \exp\left(-\frac{1}{2}\Delta \mathbf{w}^{\mathsf{T}} \mathbf{A} \Delta \mathbf{w}\right) d\mathbf{w}$$

which has no solution!

We need another approximation...

# Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$

If we assume that p(w|y) is narrow (this depends on A) then we can introduce a *linear approximation* of f(w; x) at  $w_{MAP}$ :

$$f(\mathbf{w};\mathbf{x}) \approx f(\mathbf{w}_{MAP};\mathbf{x}) + \mathbf{g}^{T} \Delta \mathbf{w}$$

where  $\mathbf{g} = \nabla f(\mathbf{w}; \mathbf{x})|_{\mathbf{w}_{MAP}}$ .

By linear approximation we just mean the Taylor expansion for k = 1.

#### This leads to

$$p(\mathbf{Y}|\mathbf{y}, \mathbf{x}) \propto \int_{\mathbb{R}^{W}} \exp\left(-\frac{\beta}{2} \left(\mathbf{y} - f(\mathbf{w}_{\text{MAP}}; \mathbf{x}) - \mathbf{g}^{\mathsf{T}} \Delta \mathbf{w}\right)^{2} - \frac{1}{2} \Delta \mathbf{w}^{\mathsf{T}} \mathbf{A} \Delta \mathbf{w}\right) d\mathbf{w}$$

and this integral can be evaluated using the  $BIG\ INTEGRAL$  to give THE ANSWER...

Method 1: approximation to  $p(\mathbf{w}|\mathbf{y})$ 

Finally

$$p(\mathbf{Y}|\mathbf{y},\mathbf{x}) = \frac{1}{\sqrt{2\pi\sigma_y^2}} \exp\left(-\frac{(\mathbf{y} - f(\mathbf{w}_{\text{MAP}};\mathbf{x}))^2}{2\sigma_y^2}\right)$$

where

$$\sigma_y^2 = \frac{1}{\beta} + \mathbf{g}^T \mathbf{A}^{-1} \mathbf{g}.$$

Hooray! But what does it mean?

# Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$

This is a Gaussian density, so we can now see that p(Y|y, x) peaks at  $f(w_{MAP}; x)$ . That is, the MAP solution.

The variance  $\sigma_u^2$  can be interpreted as a measure of certainty.

- The first term of  $\sigma_y^2$  is  $1/\beta$  and corresponds to the noise.
- The second term of  $\sigma_y^2$  is  $\mathbf{g}^T \mathbf{A}^{-1} \mathbf{g}$  and corresponds to the width of  $p(\mathbf{w}|\mathbf{y}).$

Or interpreted graphically...

# Method 1: approximation to $p(\mathbf{w}|\mathbf{y})$



# Method II: Markov chain Monte Carlo (MCMC) methods

The second solution to the problem of performing integrals

$$I = \int F(\mathbf{w}) p(\mathbf{w}|\mathbf{y}) d\mathbf{w}$$

is to use  $Monte\ Carlo\$  methods. The basic approach is to make the approximation

$$I \approx \frac{1}{N} \sum_{i=1}^{N} F(\mathbf{w}_i)$$

where the  $\mathbf{w}_i$  have distribution  $p(\mathbf{w}|\mathbf{y})$ . Unfortunately, generating  $\mathbf{w}_i$  with a given distribution can be non-trivial.

### <u>MCMC methods</u>

A simple technique is to introduce a random walk, so

 $\mathbf{w}_{i+1} = \mathbf{w}_i + \varepsilon$ 

where  $\epsilon$  is zero mean spherical Gaussian and has small variance. Obviously the sequence  $w_i$  does not have the required distribution. However we can use the *Metropolis algorithm*, which does *not* accept all the steps in the random walk:

1. If  $p(\mathbf{w}_{i+1}|\mathbf{y}) > p(\mathbf{w}_i|\mathbf{y})$  then accept the step.

2. Else accept the step with probability  $\frac{p(\mathbf{w}_{i+1}|\mathbf{y})}{p(\mathbf{w}_i|\mathbf{y})}$ .

### <u>MCMC methods</u>

In practice, the Metropolis algorithm has several shortcomings, and a great deal of research exists on improved methods, see:

R. Neal, "Probabilistic inference using Markov chain Monte Carlo methods," University of Toronto, Department of Computer Science Technical Report CRG-TR-93-1, 1993.

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