7. Gaussian Processes - Regression

Motivation

- So far we have handled learning techniques that use a parametric representation of a function
- However: what if the parametric model is incorrect or hard to find?
- Gaussian Processes model functions implicitly without parameters ("data-driven")
- Idea: represent the data using a large multivariate Normal distribution
- Only requirement: all sub-sets of data points are assumed to be jointly Normal distributed



Before We Start...

 For the following derivation we first need some math:

 Question: What is the inverse of a block matrix such as this?

 $\begin{pmatrix} A & B \\ C & D \end{pmatrix}$



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- For the following derivation we first need some math:
- Question: What is the inverse of a block matrix such as this? $\begin{pmatrix} A & B \\ C & D \end{pmatrix}$

$$\begin{pmatrix} A & B \\ C & D \end{pmatrix}^{-1} = \begin{pmatrix} (M/D)^{-1} & -(M/D)^{-1}BD^{-1} \\ -D^{-1}C(M/D)^{-1} & D^{-1} + D^{-1}C(M/D)^{-1}BD^{-1} \end{pmatrix}$$

• We have proven the matrix inversion lemma:

$$(A + BCD)^{-1} = A^{-1} - A^{-1}B(C^{-1} + DA^{-1}B)^{-1}DA^{-1}$$

Kernelized Regression (Rep.)

We can use that to compute:

$$\mathbf{w} = (\lambda I_D + \Phi^T \Phi)^{-1} \Phi^T \mathbf{t} = \Phi^T (\lambda I_N + \Phi \Phi^T)^{-1} \mathbf{t}$$

by defining: $K = \Phi \Phi^T$ $\mathbf{a} = (\lambda I_N + K)^{-1} \mathbf{t}$

we get:

$$y(\mathbf{x}) = \phi(\mathbf{x})^T \mathbf{w} = \phi(\mathbf{x})^T \Phi^T \mathbf{a}$$
$$= \mathbf{k}(\mathbf{x})^T (K + \lambda I_N)^{-1} \mathbf{t}$$

This means that the predicted output is a **linear combination** of the training outputs, where the coefficients depend on the similarities to the training input.



Gaussian Marginals and Conditionals

We can also use the result for the following:

Assume we have two variables \mathbf{x}_a and \mathbf{x}_b that are jointly Gaussian distributed, i.e. $\mathcal{N}(\mathbf{x} \mid \boldsymbol{\mu}, \boldsymbol{\Sigma})$

with

$$\mathbf{x} = \begin{pmatrix} \mathbf{x}_a \\ \mathbf{x}_b \end{pmatrix} \qquad \boldsymbol{\mu} = \begin{pmatrix} \boldsymbol{\mu}_a \\ \boldsymbol{\mu}_b \end{pmatrix} \qquad \boldsymbol{\Sigma} = \begin{pmatrix} \boldsymbol{\Sigma}_{aa} & \boldsymbol{\Sigma}_{ab} \\ \boldsymbol{\Sigma}_{ba} & \boldsymbol{\Sigma}_{bb} \end{pmatrix}$$

Question: how can we compute the conditional?

$$p(\mathbf{x}_a \mid \mathbf{x}_b)$$



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The conditional distribution is $p(\mathbf{x}_a \mid \mathbf{x}_b) = \mathcal{N}(\mathbf{x} \mid \boldsymbol{\mu}_{a|b}, \boldsymbol{\Sigma}_{a|b})$

where
$$\mu_{a|b} = \mu_a + \Sigma_{ab}\Sigma_{bb}^{-1}(\mathbf{x}_b - \mu_b)$$

and
$$\Sigma_{a|b}=\Sigma_{aa}-\Sigma_{ab}\Sigma_{bb}^{-1}\Sigma_{ba}$$
 "Schur Complement"

The marginal is $p(\mathbf{x}_a) = \mathcal{N}(\mathbf{x}_a \mid \boldsymbol{\mu}_a, \boldsymbol{\Sigma}_{aa})$





Definition

Definition: A **Gaussian process** is a collection of random variables, any finite number of which have a joint Gaussian distribution.

The number of random variables can be infinite!

This means: a GP is a Gaussian distribution over **functions**!

To specify a GP we need:

mean function: $m(\mathbf{x}) = \mathbb{E}[y(\mathbf{x})]$

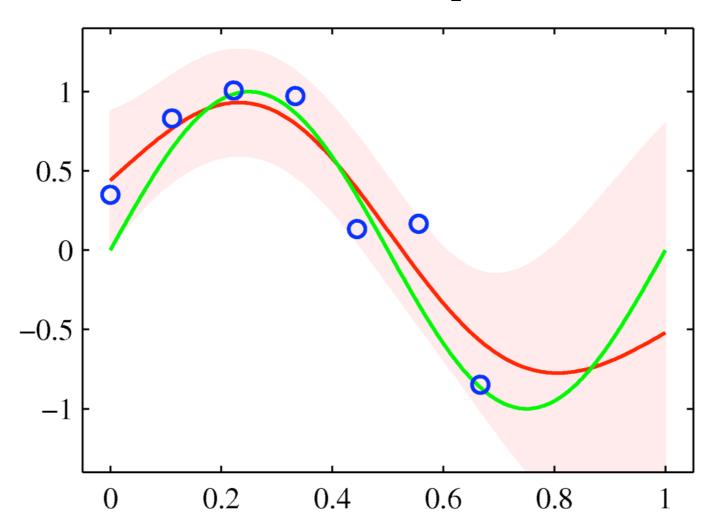
covariance function:

$$k(\mathbf{x}_1, \mathbf{x}_2) = \mathbb{E}[y(\mathbf{x}_1) - m(\mathbf{x}_1)y(\mathbf{x}_2) - m(\mathbf{x}_2)]$$





Example



- green line: sinusoidal data source
- blue circles: data points with Gaussian noise
- red line: mean function of the Gaussian process
- shaded red area: 2σ confidence interval



How Can We Handle Infinity?

Idea: split the (infinite) number of random variables into a finite and an infinite subset.

$$\mathbf{x} = \begin{pmatrix} \mathbf{x}_f \\ \mathbf{x}_i \end{pmatrix} \sim \mathcal{N} \left(\begin{pmatrix} \boldsymbol{\mu}_f \\ \boldsymbol{\mu}_i \end{pmatrix}, \begin{pmatrix} \boldsymbol{\Sigma}_f & \boldsymbol{\Sigma}_{fi} \\ \boldsymbol{\Sigma}_{fi}^T & \boldsymbol{\Sigma}_i \end{pmatrix} \right)$$

finite part

infinite part

From the marginalization property we get:

$$p(\mathbf{x}_f) = \int p(\mathbf{x}_f, \mathbf{x}_i) d\mathbf{x}_i = \mathcal{N}(\mathbf{x}_f \mid \boldsymbol{\mu}_f, \boldsymbol{\Sigma}_f)$$

This means we can use finite vectors.



A Simple Example

In Bayesian linear regression, we had $y(\mathbf{x}) = \phi(\mathbf{x})^T \mathbf{w}$ with prior probability $\mathbf{w} \sim \mathcal{N}(\mathbf{0}, \Sigma_p)$. This means:

$$\mathbb{E}[y(\mathbf{x})] = \phi(\mathbf{x})^T \mathbb{E}[\mathbf{w}] = \mathbf{0}$$

$$\mathbb{E}[y(\mathbf{x}_1)y(\mathbf{x}_2))] = \phi(\mathbf{x}_1)^T \mathbb{E}[\mathbf{w}\mathbf{w}^T]\phi(\mathbf{x}_2) = \phi(\mathbf{x}_1)^T \Sigma_p \phi(\mathbf{x}_2)$$

Any number of function values $y(\mathbf{x}_1), \dots, y(\mathbf{x}_N)$ is jointly Gaussian with zero mean.

The covariance function of this process is

$$k(\mathbf{x}_1, \mathbf{x}_2) = \phi(\mathbf{x}_1)^T \Sigma_p \phi(\mathbf{x}_2)$$

In general, any valid kernel function can be used.





The Covariance Function

The most used covariance function (kernel) is:

$$k(\mathbf{x}_p,\mathbf{x}_q)=\sigma_f^2\exp(-\frac{1}{2l^2}(\mathbf{x}_p-\mathbf{x}_q)^2)+\sigma_n^2\delta_{pq}$$
 signal variance
$$\text{length scale} \qquad \text{noise variance}$$

It is known as "squared exponential", "radial basis function" or "Gaussian kernel".

Other possibilities exist, e.g. the exponential kernel:

 $k(\mathbf{x}_p, \mathbf{x}_q) = \exp(-\theta |\mathbf{x}_p - \mathbf{x}_q|)$

This is used in the "Ornstein-Uhlenbeck" process.





Sampling from a GP

Just as we can sample from a Gaussian distribution, we can also generate samples from a GP. **Every sample will then be a function!** Process:

- 1. Choose a number of input points $\mathbf{x}_1^*, \dots, \mathbf{x}_M^*$
- 2.Compute the covariance matrix K where

$$K_{ij} = k(\mathbf{x}_i^*, \mathbf{x}_j^*)$$

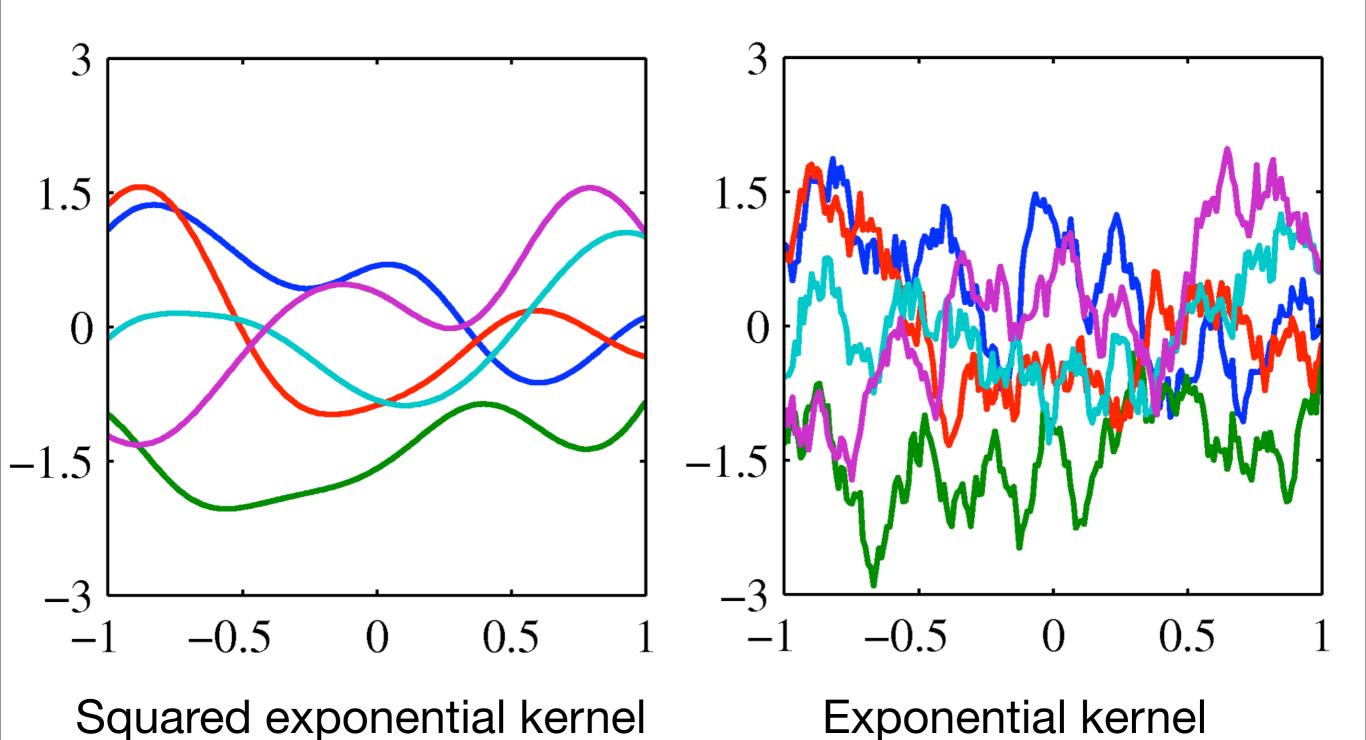
3. Generate a random Gaussian vector from

$$\mathbf{y}_* \sim \mathcal{N}(\mathbf{0}, K)$$

4. Plot the values $\mathbf{x}_1^*, \dots, \mathbf{x}_M^*$ versus y_1^*, \dots, y_M^*



Sampling from a GP







Prediction with a Gaussian Process

Most often we are more interested in predicting new function values for given input data.

We have:

training data
$$\mathbf{x}_1, \dots, \mathbf{x}_N$$
 t_1, \dots, t_N test input $\mathbf{x}_1^*, \dots, \mathbf{x}_M^*$

And we want test outputs y_1^*, \dots, y_M^*

The joint probability is

$$\begin{pmatrix} \mathbf{y} \\ \mathbf{y}_* \end{pmatrix} \sim \mathcal{N} \left(\mathbf{0}, \begin{pmatrix} K(X, X) & K(X, X_*) \\ K(X_*, X) & K(X_*, X_*) \end{pmatrix} \right)$$

and we need to compute $p(\mathbf{y}^* \mid \mathbf{x}^*, X, \mathbf{y})$.





Prediction with a Gaussian Process

In the case of only one test point x^* we have

$$K(X, \mathbf{x}^*) = \begin{pmatrix} k(\mathbf{x}_1, \mathbf{x}_*) \\ \vdots \\ k(\mathbf{x}_N, \mathbf{x}_*) \end{pmatrix} = \mathbf{k}_*$$

Now we compute the conditional distribution

$$p(y^* \mid \mathbf{x}^*, X, \mathbf{y}) = \mathcal{N}(y_* \mid \mu_*, \Sigma_*)$$

where

$$\mu_* = \mathbf{k}_*^T K^{-1} \mathbf{y}$$

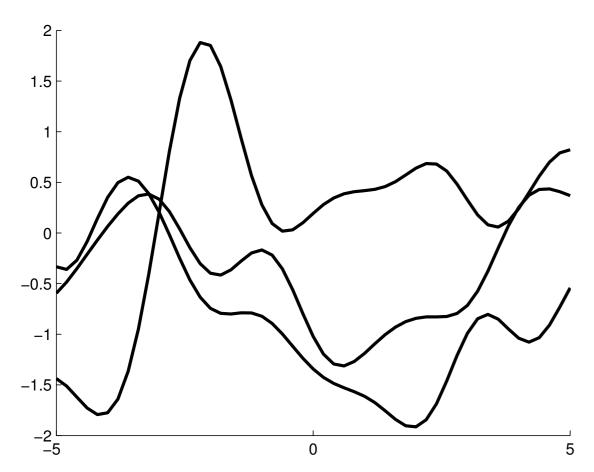
$$\Sigma_* = k(\mathbf{x}_*, \mathbf{x}_*) - \mathbf{k}_*^T K^{-1} \mathbf{k}_*$$

This defines the predictive distribution.

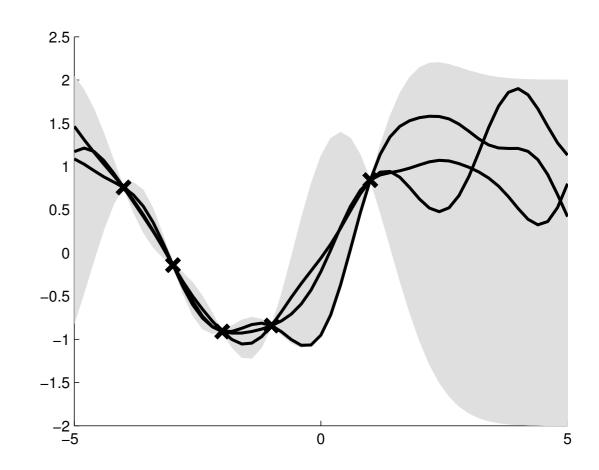




Example



Functions sampled from a Gaussian Process prior

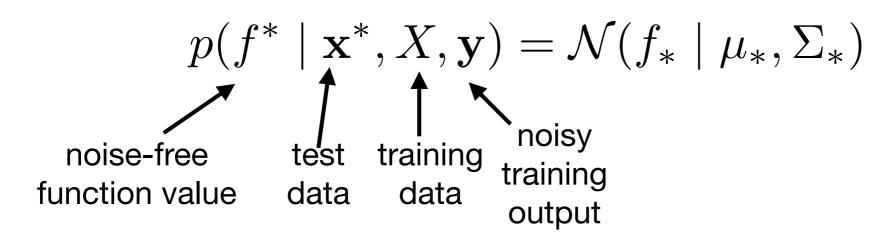


Functions sampled from the predictive distribution

The predictive distribution is itself a Gaussian process. It represents the posterior after observing the data. The covariance is low in the vicinity of data points.

Prediction With Noisy Observations

- Assume we have observations with noise
- We can model this with a latent function f
- Thus: y is a noisy variant of f
- The formulation only differs slightly:

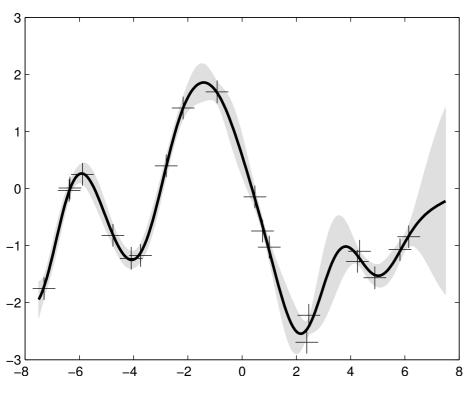


We assume that the noise is Gaussian:

$$p(\mathbf{y} \mid \mathbf{f}) = \mathcal{N}(\mathbf{f}, \sigma_y^2 I)$$

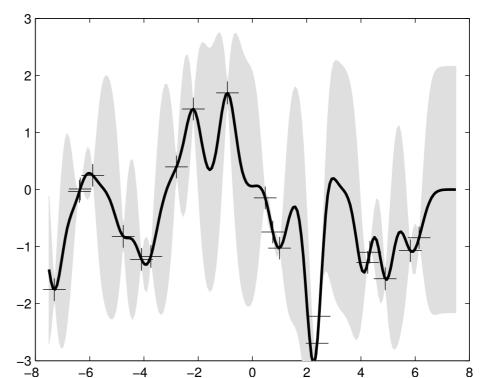


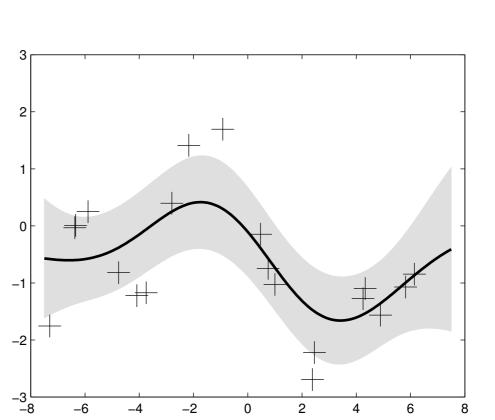
Varying the Hyperparameters



$$l = \sigma_f = 1, \quad \sigma_n = 0.1$$

- 20 data samples
- GP prediction with different kernel hyper parameters





$$l = 0.3,$$

$$\sigma_f = 1.08,$$

$$\sigma_n = 0.0005$$

$$l=3$$

$$\sigma_f = 1.16$$

$$\sigma_n = 0.89$$

Varying the Hyperparameters

The squared exponential covariance function can be generalized to

$$k(\mathbf{x}_p, \mathbf{x}_q) = \sigma_f^2 \exp(-\frac{1}{2}(\mathbf{x}_p - \mathbf{x}_q)^T M(\mathbf{x}_p - \mathbf{x}_q)) + \sigma_n^2 \delta_{pq}$$

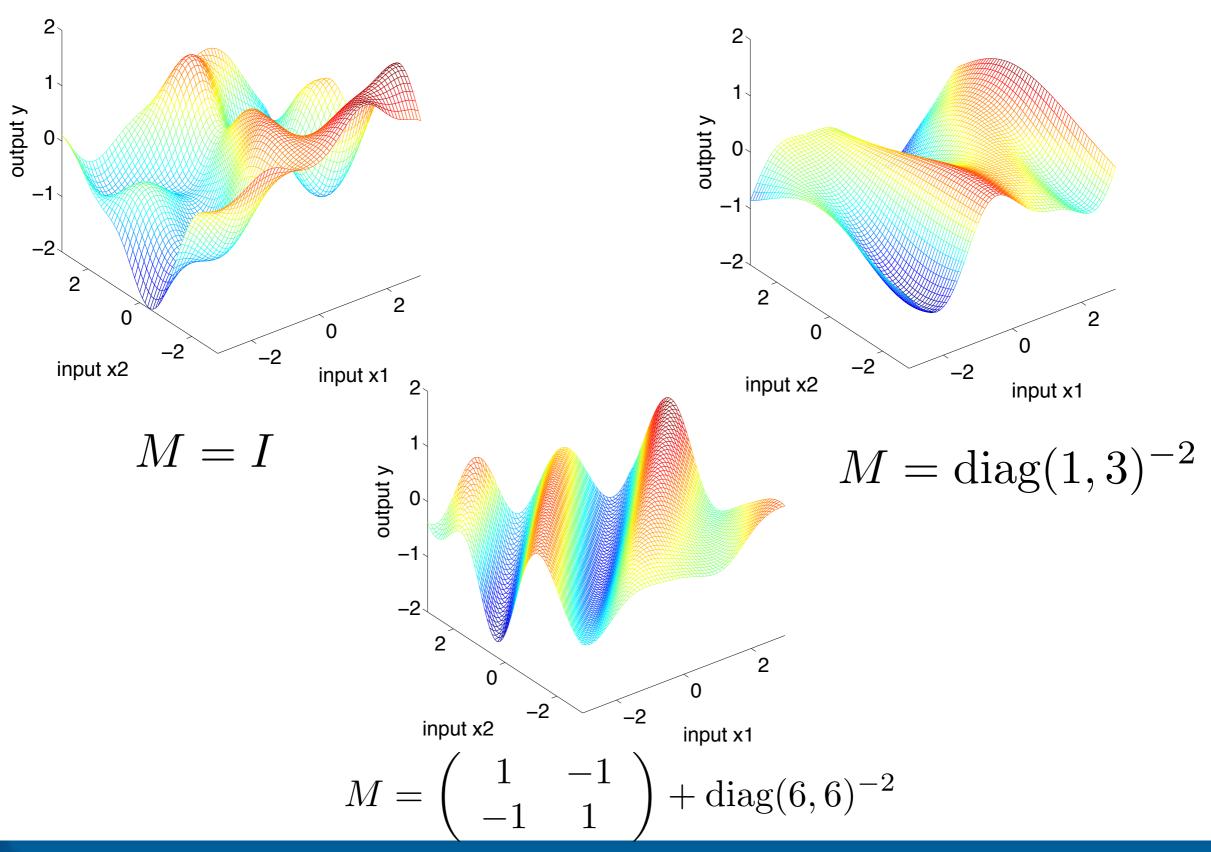
where M can be:

- $M = l^{-2}I$: this is equal to the above case
- $M = \operatorname{diag}(l_1, \dots, l_D)^{-2}$: every feature dimension has its own length scale parameter
- $M = \Lambda \Lambda^T + \mathrm{diag}(l_1, \dots, l_D)^{-2}$: here Λ has less than D columns





Varying the Hyperparameters



Implementation

Algorithm 1: GP regression

 $\mathtt{var}[f_*] \leftarrow k(\mathbf{x}_*, \mathbf{x}_*) - \mathbf{v}^T \mathbf{v}$

Data: training data (X, \mathbf{y}) , test data \mathbf{x}_*

Input: Hyper parameters σ_f^2 , l, σ_n^2

$$K_{ij} \leftarrow k(\mathbf{x}_i, \mathbf{x}_j)$$
 — $L \leftarrow \text{cholesky}(K + \sigma_y^2 I)$ $\boldsymbol{\alpha} \leftarrow L^T \backslash (L \backslash \mathbf{y})$ —

Training Phase

 $\mathbb{E}[f_*] \leftarrow \mathbf{k}_*^T \boldsymbol{\alpha}$ $\mathbf{v} \leftarrow L \backslash \mathbf{k}_*$

Test Phase

- $\log p(\mathbf{y} \mid X) \leftarrow -\frac{1}{2} \mathbf{y}^T \boldsymbol{\alpha} \sum_i \log L_{ii} \frac{N}{2} \log(2\pi)$
- Cholesky decomposition is numerically stable
- Can be used to compute inverse efficiently

To find optimal hyper parameters we need the marginal likelihood:

$$p(\mathbf{y} \mid X) = \int p(\mathbf{y} \mid \mathbf{f}, X) p(\mathbf{f} \mid X) d\mathbf{f}$$

$$p(\mathbf{y} \mid \mathbf{f}) = \prod_{i} \mathcal{N}(y_i \mid f_i, \sigma_y^2) \qquad p(\mathbf{f} \mid X) = \mathcal{N}(\mathbf{f} \mid 0, K)$$

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

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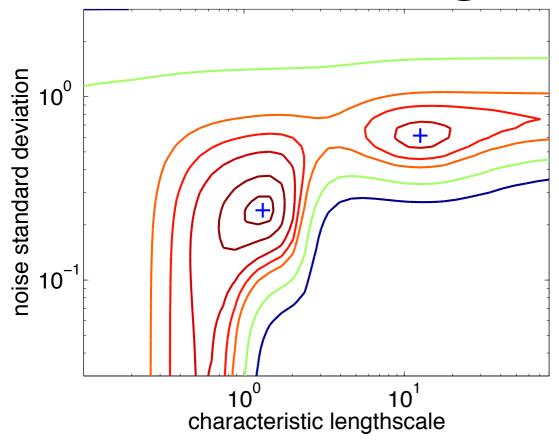
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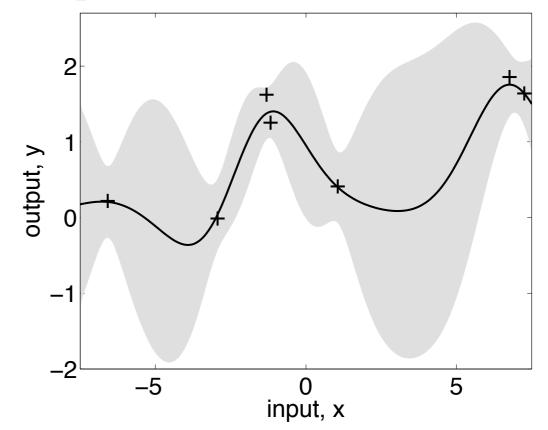
We take the logarithm, compute the derivative and set it to θ . This is the **training** step.

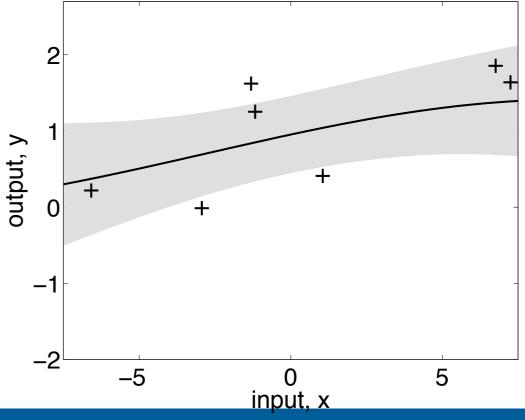




The log marginal likelihood is not necessarily concave, i.e. it can have local maxima.

The local maxima can correspond to sub-optimal solutions.





Automatic Relevance Determination

- We have seen how the covariance function can be generalized using a matrix M
- If M is diagonal this results in the kernel function

$$k(\mathbf{x}, \mathbf{x}') = \sigma_f \exp\left(\frac{1}{2} \sum_{i=1}^{D} \eta_i (x_i - x_i')^2\right)$$

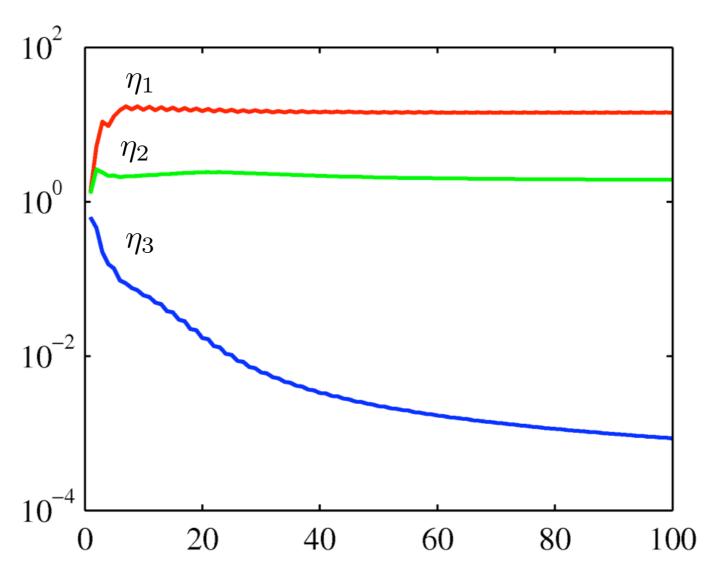
- We can interpret the η_i as weights for each feature dimension
- Thus, if the length scale $l_i = 1/\eta_i$ of an input dimension is large, the input is less relevant
- During training this is done automatically





Automatic Relevance Determination

3-dimensional data, parameters η_1 η_2 η_3 as they evolve during training



During the optimization process to learn the hyper-parameters, the reciprocal length scale for one parameter decreases, i.e.:

This hyper parameter is not very relevant!





Gaussian Processes - Classification

Gaussian Processes For Classification

In regression we have $y \in \mathbb{R}$, in binary classification we have $y \in \{-1, 1\}$

To use a GP for classification, we can apply a **sigmoid** function to the posterior obtained from the GP and compute the class probability as:

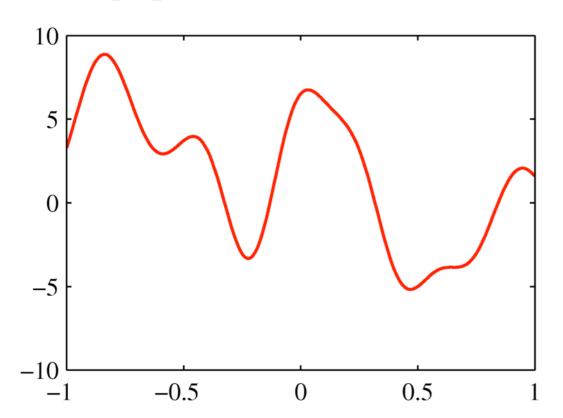
$$p(y = +1 \mid \mathbf{x}) = \sigma(f(\mathbf{x}))$$

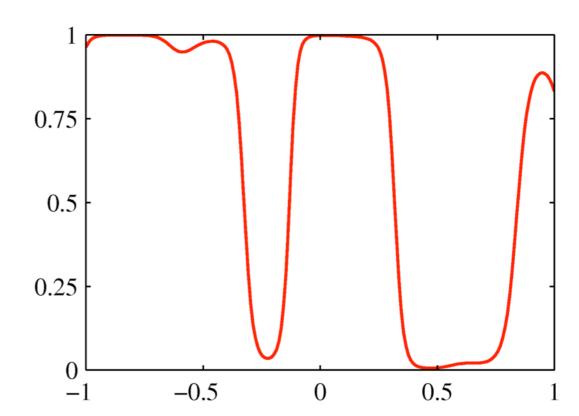
If the sigmoid function is symmetric: $\sigma(-z) = 1 - \sigma(z)$ then we have $p(y \mid \mathbf{x}) = \sigma(yf(\mathbf{x}))$.

A typical type of sigmoid function is the logistic sigmoid: $\sigma(z) = \frac{1}{1 + \exp(-z)}$



Application of the Sigmoid Function





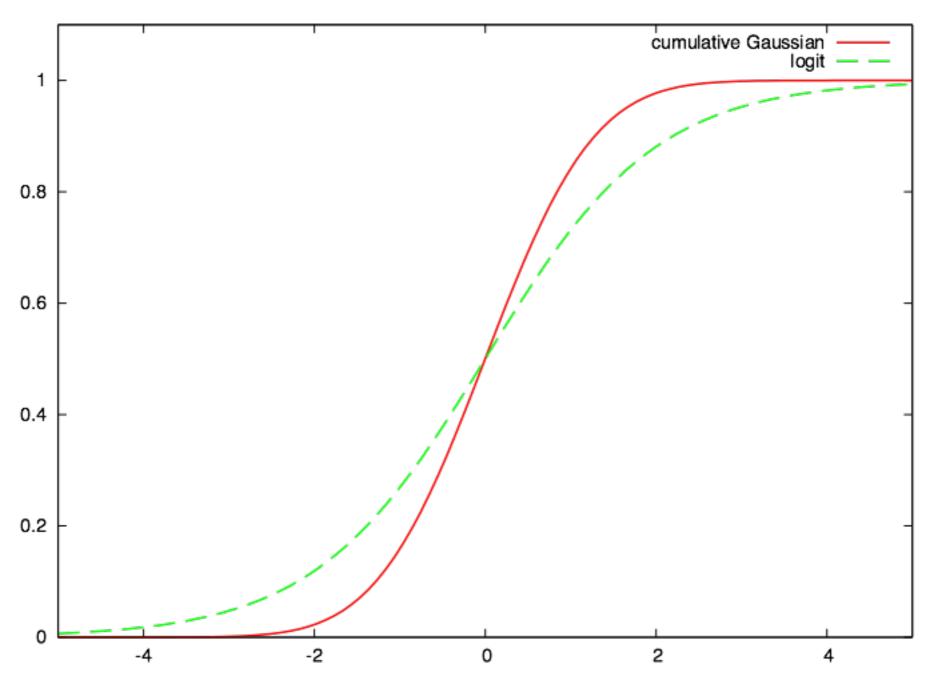
Function sampled from a Gaussian Process

Sigmoid function applied to the GP function

Another symmetric sigmoid function is the cumulative Gaussian:

$$\Phi(z) = \int_{-\infty}^{z} \mathcal{N}(x \mid 0, 1) dx$$

Visualization of Sigmoid Functions



The cumulative Gaussian is slightly steeper than the logistic sigmoid

The Latent Variables

In regression, we directly estimated f as

$$f(\mathbf{x}) \sim \mathcal{GP}(m(\mathbf{x}), k(\mathbf{x}, \mathbf{x}'))$$

and values of *f* where observed in the training data. Now only labels +1 or -1 are observed and *f* is treated as a set of **latent variables.**

A major advantage of the Gaussian process classifier over other methods is that it **marginalizes** over all latent functions rather than maximizing some model parameters.



Class Prediction with a GP

The aim is to compute the predictive distribution

$$p(y_* = +1 \mid X, \mathbf{y}, \mathbf{x}_*) = \int p(y_* \mid f_*) p(f_* \mid X, \mathbf{y}, \mathbf{x}_*) df_*$$



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we marginalize over the latent variables from the training data:

$$p(f_* \mid X, \mathbf{y}, \mathbf{x}_*) = \int p(f_* \mid X, \mathbf{x}_*, \mathbf{f}) p(\mathbf{f} \mid X, \mathbf{y}) d\mathbf{f}$$

predictive distribution of the latent variable (from regression)

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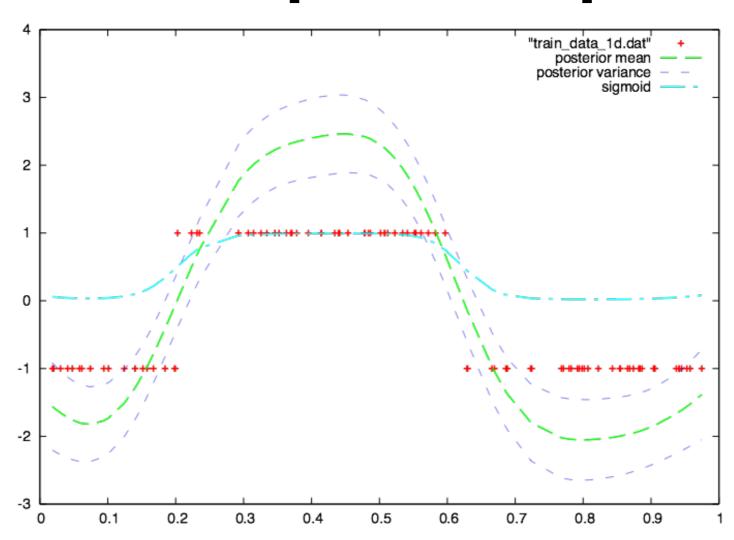
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we need the posterior over the latent variables:

(sigmoid)
$$p(\mathbf{f}\mid X,\mathbf{y}) = \frac{p(\mathbf{y}\mid \mathbf{f})p(\mathbf{f}\mid X)}{p(\mathbf{y}\mid X)}$$
 prior normalizer

A Simple Example



- Red: Two-class training data
- Green: mean function of $p(\mathbf{f} \mid X, \mathbf{y})$
- Light blue: sigmoid of the mean function

But There Is A Problem...

$$p(\mathbf{f} \mid X, \mathbf{y}) = \frac{p(\mathbf{y} \mid \mathbf{f})p(\mathbf{f} \mid X)}{p(\mathbf{y} \mid X)}$$

- The likelihood term is not a Gaussian!
- This means, we can not compute the posterior in closed form.
- There are several different solutions in the literature, e.g.:
 - Laplace approximation
 - Expectation Propagation
 - Variational methods





Laplace Approximation

$$p(\mathbf{f} \mid X, \mathbf{y}) \approx q(\mathbf{f} \mid X, \mathbf{y}) = \mathcal{N}(\mathbf{f} \mid \hat{\mathbf{f}}, A^{-1})$$

where
$$\hat{\mathbf{f}} = \arg\max_{\mathbf{f}} p(\mathbf{f} \mid X, \mathbf{y})$$

and $A = -\nabla\nabla\log p(\mathbf{f} \mid X, \mathbf{y})|_{\mathbf{f} = \hat{\mathbf{f}}}$

second-order
Taylor expansion

To compute \hat{f} an iterative approach using Newton's method has to be used.

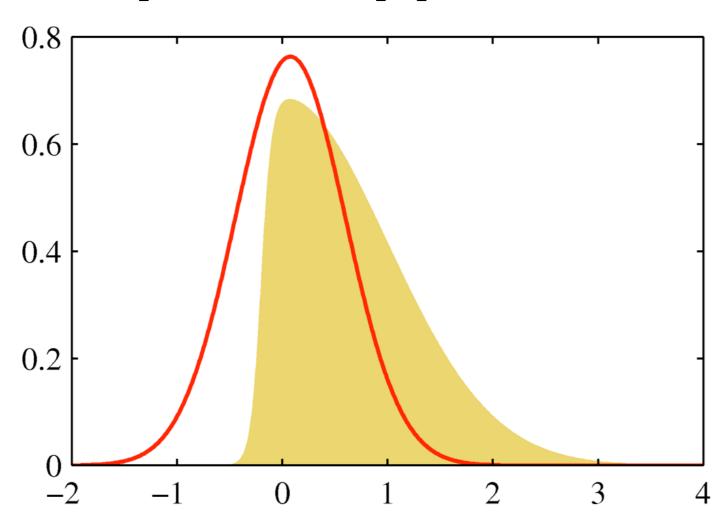
The Hessian matrix A can be computed as

$$A = K^{-1} + W$$

where $W = -\nabla\nabla \log p(\mathbf{y} \mid \mathbf{f})$ is a diagonal matrix which depends on the sigmoid function.



Laplace Approximation



- Yellow: a non-Gaussian posterior
- Red: a Gaussian approximation, the mean is the mode of the posterior, the variance is the negative second derivative at the mode

Predictions

Now that we have $p(\mathbf{f} \mid X, \mathbf{y})$ we can compute:

$$p(f_* \mid X, \mathbf{y}, \mathbf{x}_*) = \int p(f_* \mid X, \mathbf{x}_*, \mathbf{f}) p(\mathbf{f} \mid X, \mathbf{y}) d\mathbf{f}$$

From the regression case we have:

$$p(f_* \mid X, \mathbf{x}_*, \mathbf{f}) = \mathcal{N}(f_* \mid \mu_*, \Sigma_*)$$
 where $\mu_* = \mathbf{k}_*^T K^{-1} \mathbf{f}$
$$\Sigma_* = k(\mathbf{x}_*, \mathbf{x}_*) - \mathbf{k}_*^T K^{-1} \mathbf{k}_*$$

Linear in f

This reminds us of a property of Gaussians that we saw earlier!



Gaussian Properties (Rep.)

If we are given this:

$$p(\mathbf{x}) = \mathcal{N}(\mathbf{x} \mid \mu, \Sigma_1)$$

II.
$$p(\mathbf{y} \mid \mathbf{x}) = \mathcal{N}(\mathbf{y} \mid A\mathbf{x} + \mathbf{b}, \Sigma_2)$$

Then it follows (properties of Gaussians):

III.
$$p(\mathbf{y}) = \mathcal{N}(\mathbf{y} \mid A\mu + \mathbf{b}, \Sigma_2 + A\Sigma_1 A^T)$$

IV.
$$p(\mathbf{x} \mid \mathbf{y}) = \mathcal{N}(\mathbf{x} \mid \Sigma(A^T \Sigma_2^{-1} (\mathbf{y} - \mathbf{b}) + \Sigma_1^{-1} \mathbf{y}), \Sigma)$$

where

$$\Sigma = (\Sigma_1^{-1} + A^T \Sigma_s^{-1} A)^{-1}$$

Applying this to Laplace

$$\mathbb{E}[f_* \mid X, \mathbf{y}, \mathbf{x}_*] = \mathbf{k}(\mathbf{x}_*)^T K^{-1} \hat{\mathbf{f}}$$

$$\mathbb{V}[f_* \mid X, \mathbf{y}, \mathbf{x}_*] = k(\mathbf{x}_*, \mathbf{x}_*) - \mathbf{k}_*^T (K + W^{-1})^{-1} \mathbf{k}_*$$

It remains to compute

$$p(y_* = +1 \mid X, \mathbf{y}, \mathbf{x}_*) = \int p(y_* \mid f_*) p(f_* \mid X, \mathbf{y}, \mathbf{x}_*) df_*$$

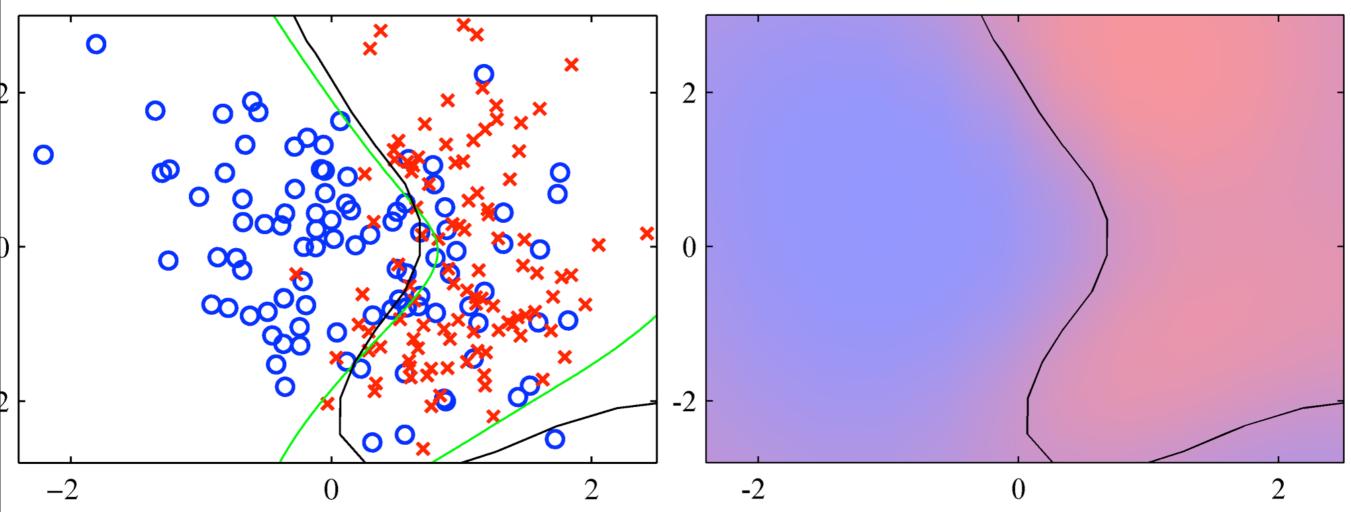
Depending on the kind of sigmoid function we

- can compute this in closed form (cumulative Gaussian sigmoid)
- have to use sampling methods or analytical approximations (logistic sigmoid)





A Simple Example



- Two-class problem (training data in red and blue)
- Green line: optimal decision boundary
- Black line: GP classifier decision boundary
- Right: posterior probability



Properties of the GP Classifier

- The biggest benefit of the GPC is that it gives a "well-calibrated" estimate of the uncertainty:
 - The prediction is based on marginalization over all possible models, not on a particular one
 - Other techniques use the distance to the decision boundary, which is often not correct
- The biggest disadvantage of the GPC is its memory and run time complexity
- However: There are sparse versions of the GPC, such as the Informative Vector Machine (IVM), that uses a similar idea than the SVM



Summary

- Gaussian Processes are Normal distributions over functions
- To specify a GP we need a covariance function (kernel) and a mean function
- For regression we can compute the predictive distribution in closed form
- For classification, we use a sigmoid and have to approximate the latent posterior
- More on Gaussian Processes:
 http://videolectures.net/epsrcws08_rasmussen_lgp/

