



MATH6221, Topic in Numerical Analysis

Computational Inverse Problems

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Outline

- 1 Fundamentals of Bayesian inversion
- 2 Markov Chain Monte Carlo
- 3 Approximate inference



finite-dimensional inverse problem

$$F(f) = g,$$

- f : the unknown signal/image, linear / nonlinear
- g : the noisy data (g^δ)
- $F : \mathbb{R}^m \mapsto \mathbb{R}^n$: forward map, ill-conditioned



Motivation

- variational regularization $\Leftrightarrow$ optimization problem, e.g.,

$$\|F(f) - g\|^2 + \alpha \|f\|^2 / \mathcal{R}(f)$$

$\Rightarrow$ Tikhonov minimizer

- Question: How plausible is the Tikhonov minimizer f_α^δ ?
How effective is one penalty compared to another ? ...
- $\Rightarrow$ tools for assessing the reliability of the inverse solution
 - Bayesian inference
 - interval analysis ...
 - sensitivity analysis ...
 - ...



Bayesian inversion

starting point: Bayes' formula, i.e., for two random variables f and g the conditional probability density of f given g is given by

$$p(f|g) = \frac{p(g|f)p(f)}{p(g)}$$

- $p(g|f)$: **likelihood function** — building block I
the probability density of the data g given the unknown f
incorporate information in g (physics / noise statistics of g)
- $p(f)$: **prior distribution** — building block II
the probability density of f (regardless of g)
encode a prior knowledge/information (before collecting data)
- $p(g) = \int p(g|f)p(f)df$: normalizing constant



- the unnormalized posteriori $p(f, g)$ defined by

$$p(f, g) = p(g|f)p(f),$$

and often write

$$p(f|g) \propto p(f, g)$$

the posteriori $p(f|g)$ up to a normalizing constant $p(g)$

- the normalizing constant $p(g)$ is *important* for model comparison
...



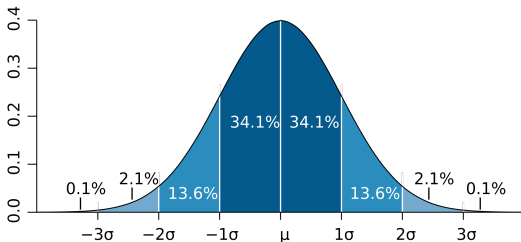
posterior $p(f|g)$: the distr. of parameter f given the data g

⇒ Bayesian solution of the inverse problem

- $p(f|g)$ holds the full inform. about the inverse problem
- there are many plausible solutions to the inverse problem
- ⇒ calibrating the uncertainty of specific inverse solution.



Gaussian distribution ...



- the mean μ of f under $p(f|g) \Rightarrow$ representative solution
- the variance σ of f under $p(f|g) \Rightarrow$ reliability/uncertainties

all these involve (possibly very high-dimensional) integrals



Building block I: the construction of likelihood function

likelihood function $p(g|f) \Leftarrow$ the noise statistics, the forward model

- all sources of errors: modeling error, discretization error, measurement error ...
all are lumped into g despite the **very different nature** ...
- a careful modeling of all errors in the data g is essential for extracting useful information



The most popular model is the **additive** noise model

$$g = g^\dagger + \xi = F(f^\dagger) + \xi,$$

- given a realization $f^\dagger$ of f : g is a shift of the noise ξ by $F(f^\dagger)$
- ξ is from external source: it is assumed to be *independent* of $f \Rightarrow$

$$p(g|f) = p_\Xi(g - F(f))$$

additive iid Gaussian noise

$\xi \in \mathbb{R}^n$ is an i.i.d. **realization** of a Gaussian r.v. $N(0, \sigma^2)$,

$$\xi_i \sim N(0, \sigma^2) \Rightarrow p(g_i|f) = \frac{1}{\sqrt{2\pi\sigma^2}} e^{-\frac{1}{2\sigma^2}(g_i - F(f)_i)^2}$$

$$\text{i.i.d. noise} \Rightarrow p(g|f) = (2\pi\sigma^2)^{-\frac{n}{2}} e^{-\frac{1}{2\sigma^2}\|F(f) - g\|^2}$$

$$-\log p(g|f) = \frac{1}{2\sigma^2}\|F(f) - g\|^2 + c$$



general Gaussian case

$\xi \in \mathbb{R}^n$ is a realization of Gaussian r.v. $N(\mathbf{0}, \mathbf{C})$

$$p(g|f) = \frac{1}{(2\pi)^{\frac{n}{2}} |\det \mathbf{C}|^{\frac{1}{2}}} e^{-\frac{1}{2}(F(f)-g)^t \mathbf{C}^{-1} (F(f)-g)}$$

with $\mathbf{C} = \sigma^2 I$, recovers the i.i.d. case
the discretization error is often colored ...

i.i.d. Laplace noise

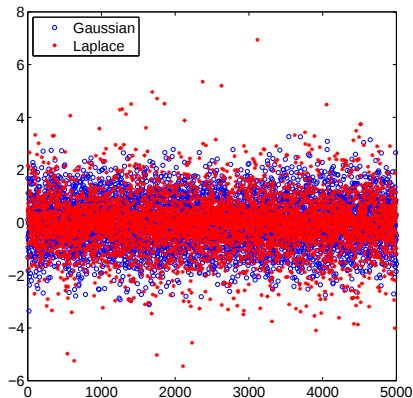
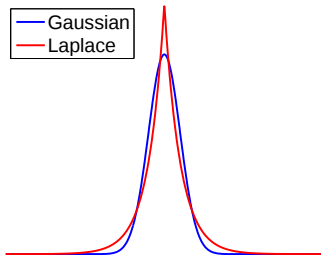
$\xi_i \sim \text{Lap}(0, \sigma)$, $p_{\Xi}(\xi) = (2\sigma)^{-1} e^{-\sigma|\xi|}$; i.i.d. noise $\Rightarrow$

$$p(g|f) = \frac{1}{(2\sigma)^n} e^{-\sigma \|F(f)-g\|_1}$$

suitable for *impulsive* noise ...



Gaussian v.s. Laplace distribution with identical means/variance





Poisson noise

$g_i \sim \text{Pois}(g_i^\dagger)$, i.e.,

$$g_i = k, \quad \text{with } p(k) = \frac{(g_i^\dagger)^k}{k!} e^{-g_i^\dagger}, \quad k = 0, 1, \dots$$

i.i.d. assumption on the noise model

$$p(g|f) = \prod_i \frac{e^{-(F(f))_i} (F(f))_i^{g_i}}{g_i!}$$

suitable in photon-counting: positron emission tomography

all these analytical models are approximations, but some are better than others ...



Building block II: the construction of prior

The prior $p(f)$ encodes the prior knowledge about the sought-for solution f in a **probabilistic** manner.

- There is no universal way to construct prior, and it is *subjective*, dependent on personal experiences/knowledge
- the prior knowledge: expert opinion, historical investigations, statistical studies and anatomical knowledge etc.
- Since inverse problems are ill-posed: careful incorporation of all prior knowledge is important in any inversion technique
- the prior plays the role of regularization in a stochastic setting



Example: one-dimensional deblurring problem

$$g(t) = \int_0^1 k(s, t) f(s) ds \quad t \in [0, 1]$$

- grid: $0 = s_0 < s_1 < \dots < s_n = 1, 0 = t_0 < t_1 < \dots < t_n = 1$
- proper discretization: collocation, Galerkin
- discrete problem

$$g = Af$$



smoothness prior: the solution is *smooth*, e.g.,

$$f_i \approx f_{i-1}$$

too rigid $\Rightarrow$ include uncertainty in the belief by

$$f_i = f_{i-1} + \xi_i$$

with $\xi_i \sim N(0, \gamma^2)$

$$Lf = \xi$$

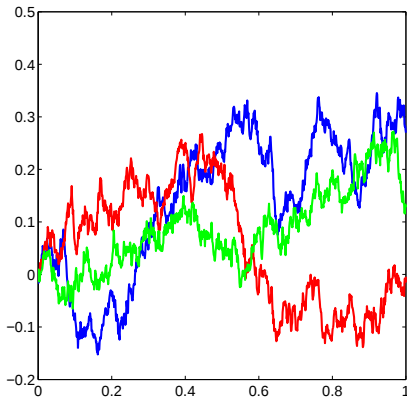
with L given by

$$L = \begin{bmatrix} -1 & 1 & & & & \\ & -1 & 1 & & & \\ & & & \ddots & & \\ & & & & \ddots & \\ & & & & & -1 & 1 \end{bmatrix}$$

first-order backward difference (up to a step size)



- prior $p(f) \propto e^{-\frac{1}{2\gamma^2} \|Lf\|^2}$
- the prior is not normalizable (improper), likelihood $p(g|f)$ hopefully fixes the problem





smoothness prior: solution is *smooth*

$$f_i = (f_{i-1} + f_{i+1})/2$$

value is close to its neighbors, but no uncertainties $\Rightarrow$ innovations

$$f_i = (f_{i-1} + f_{i+1})/2 + \xi_i$$

with $\xi_i \sim N(0, \gamma^2)$, i.e.,

$$Lf = \xi$$

with

$$L = \frac{1}{2} \begin{bmatrix} -1 & 2 & -1 & & & \\ & -1 & 2 & -1 & & \\ & & \ddots & \ddots & \ddots & \\ & & & -1 & 2 & 1 \end{bmatrix}$$

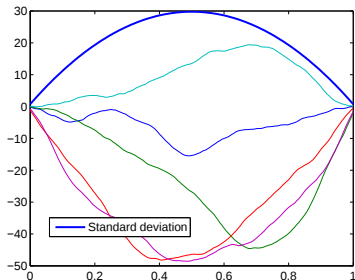
... discrete Laplacian ... (second-order derivative is small ...)



prior distribution

$$p(f) \propto e^{-\frac{1}{2\gamma^2} \|Lf\|^2}$$

- this is not a probability density function (not normalizable ...)
- cause: no proper boundary conditions for the Laplacian
- remedy: specify proper boundary condition / let likelihood fix it





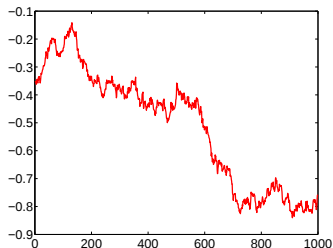
nonsmooth prior: locally smooth, but large jumps at *unknown* location

$$f_i = f_{i-1} + \xi_i$$

ξ_i mostly small with occasional large jumps $\Rightarrow \xi_i \sim \text{Lap}(0, \gamma)$?

$$Lf = \xi,$$

total variation prior distribution $p(f) \propto e^{-\gamma \|Lf\|_1}$



the samples not exactly piecewise constant like TV regularization ...



all can be regarded as special cases of **Markov random field**

$$p(f) \propto e^{-\lambda \psi(f)},$$

- $\psi(f)$: potential function dictates the interaction energy between the components of the random field f
- λ is a scale parameter, determining the strength of the local/global interactions.

It plays the role of a regularization parameter in regularization theory, and its automated determination is important.



distinct features

- $p(f|g)$ is a **probability distribution**, and is an ensemble of solutions consistent with g (to various extent)

$$\mu = \int f p(f|g) df,$$

$$C = \int (f - \mu)(f - \mu)^t p(f|g) df.$$

- elucidate the crucial role of proper **statistical modeling** in designing regularization formulations for practical problems.
- it provides a flexible regularization, **partially** resolving the nontrivial issue of choosing a regularization parameter.



hierarchical modeling: hyperparameters and beyond

likelihood $p(g|f)$ and prior $p(f)$ may contain unknown para, e.g.,

$$p(g|f) = p(g|f, \tau) \quad \text{and} \quad p(f) = p(f|\lambda)$$

- τ, λ : precision (inverse variance) and the scale parameter
- These parameters are generically known as hyperparameters
- Hierarchical Bayesian modeling provides an elegant approach to choose these parameters automatically

J. Wang, N. Zabaras. Inverse Problems 21(1), 183–206, 2005



hierarchical Bayesian modeling

- view λ and τ as random variables with their own priors
- determine them from the data g
- convenient choice: **conjugate** distribution

For both λ and τ , the conjugate distribution is a Gamma distribution:

$$p_{\Lambda}(\lambda) = G(\lambda; a_0, b_0) = \frac{b_0^{a_0}}{\Gamma(a_0)} \lambda^{a_0-1} e^{-b_0 \lambda},$$

$$p_{\Upsilon}(\tau) = G(\tau; a_1, b_1) = \frac{b_1^{a_1}}{\Gamma(a_1)} \tau^{a_1-1} e^{-b_1 \tau}.$$

- (a_0, b_0) and (a_1, b_1) determine the range of λ and τ
- noninformative prior is often adopted: $(a_0, b_0) \approx (1, 0)$

posterior distribution $p(f, \lambda, \tau | g)$

$$p(f, \lambda, \tau | g) \propto p(g | f, \tau) p(f | \lambda) p_{\Lambda}(\lambda) p_{\Upsilon}(\tau)$$



example: Gaussian noise model + Gaussian prior

$$p(g|f, \tau) \propto \tau^{\frac{n}{2}} e^{-\frac{\tau}{2} \|F(f) - g\|^2},$$

$$p(f|\lambda) \propto \lambda^{\frac{m}{2}} e^{-\frac{\lambda}{2} \|f\|^2}$$

the posterior distribution (for fixed λ and τ)

$$p(f|g) \propto e^{-(\frac{\tau}{2} \|F(f) - g\|^2 + \frac{\lambda}{2} \|f\|^2)}$$

a customary approach: maximum a posteriori estimate

$$f_{\text{map}} = \arg \max_f p(f|g)$$

the *most probable* point of the posterior density



maximum a posteriori estimate $f_{\text{map}} \Rightarrow$

$$\begin{aligned} f_{\text{map}} &= \arg \max_f p(f|g) \\ &= \arg \min_f \left\{ \frac{\tau}{2} \|F(f) - g\|^2 + \frac{\lambda}{2} \|f\|^2 \right\} \end{aligned}$$

the functional in the curly bracket is

$$\frac{1}{2} \|F(f) - g\|^2 + \frac{\lambda \tau^{-1}}{2} \|f\|^2$$

Tikhonov regularization + smoothness constraint, with $\alpha = \lambda \tau^{-1}$

Message: A Tikhonov minimizer is an MAP estimate of some Bayesian formulation.



What if the parameters λ and τ are *unknown*?

⇒ hierarchical model conjugate prior on λ and τ

$$p(g|f, \tau) \propto \tau^{\frac{n}{2}} e^{-\frac{\tau}{2} \|F(f) - g\|^2}$$

$$p(f|\lambda) \propto \lambda^{\frac{m}{2}} e^{-\frac{\lambda}{2} \|f\|^2}$$

$$p(\tau) \propto \tau^{a_1-1} e^{-b_1 \tau}$$

$$p(\lambda) \propto \lambda^{a_0-1} e^{-b_0 \lambda}$$

with $(a_1, b_1), (a_0, b_0) \simeq (0, 0)$

⇒ posterior distr.

$$\begin{aligned} p(f, \lambda, \tau|g) &\propto \tau^{\frac{n}{2} + a_1 - 1} e^{-\frac{\tau}{2} \|F(f) - g\|^2} \\ &\cdot \lambda^{\frac{m}{2} + a_0 - 1} e^{-\frac{\lambda}{2} \|f\|^2} \cdot e^{-b_1 \tau} \cdot e^{-b_0 \lambda}. \end{aligned}$$



ways of exploring the posterior distribution $p(f, \lambda, \tau|g)$

- the *joint maximum a posteriori* estimate $(f, \lambda, \tau)_{\text{map}}$, i.e.,

$$\begin{aligned}(f, \lambda, \tau)_{\text{map}} &= \arg \max_{f, \lambda, \tau} p(f, \lambda, \tau|g) \\ &= \arg \min_{f, \lambda, \tau} J(f, \lambda, \tau),\end{aligned}$$

where the functional $J(f, \lambda, \tau)$ is given by

$$J(f, \lambda, \tau) = \frac{\tau}{2} \|F(f) - g\|^2 + \frac{\lambda}{2} \|f\|^2 - \tilde{a}_0 \ln \lambda + b_0 \lambda - \tilde{a}_1 \ln \tau + b_1 \tau,$$

with $\tilde{a}_0 = \frac{n}{2} + a_0 - 1$ and $\tilde{a}_1 = \frac{m}{2} + a_1 - 1$



augmented Tikhonov regularization

$$J(f, \lambda, \tau) = \frac{\tau}{2} \|F(f) - g\|^2 + \frac{\lambda}{2} \|f\|^2 - \tilde{a}_0 \ln \lambda + b_0 \lambda - \tilde{a}_1 \ln \tau + b_1 \tau.$$

- the first two terms recover Tikhonov regularization
- the rest terms automatically determine the regu. parameter.

What is beyond: MAP approach is straightforward, but without uncertainties ???



How to explore the posteriori ?

posteriori $p(f)$ lives in a high-dimensional space

⇒ not directly informative

⇒ compute summarizing statistics, e.g., mean μ and covariance C

$$\mu = \int f p(f) df \quad \text{and} \quad C = \int (f - \mu)(f - \mu)^t p(f) df.$$

very high-dim. integrals, and standard quadrature rules are inefficient

$$m = 100, 2 \text{ points/dir} \Rightarrow 2^{100} \approx 1.27 \times 10^{30} \text{ points}$$

more efficient approach

- Monte Carlo methods, especially Markov chain Monte Carlo

J. S. Liu. Monte Carlo Strategies in Scientific Computing. Springer-Verlag, New York, 2001.



Monte Carlo simulation

- draw a large set of i.i.d. samples $\{f^{(i)}\}_{i=1}^N$ from the target distribution $p(f)$
- approximate the expectation $E_p[\zeta]$ of any function $\zeta : \mathbb{R}^m \rightarrow \mathbb{R}$ by the sample mean $E_N[\zeta]$

$$E_N[\zeta] \equiv \frac{1}{N} \sum_{i=1}^N \zeta(f^{(i)}) \rightarrow E_p[\zeta] = \int \zeta(f) p(f) df \quad \text{as } N \rightarrow \infty.$$

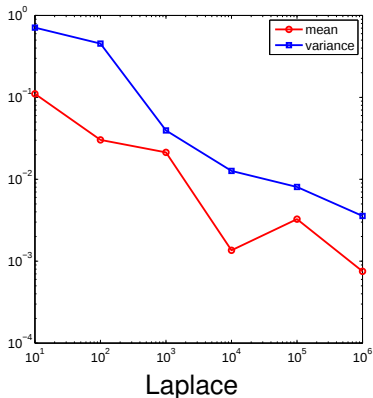
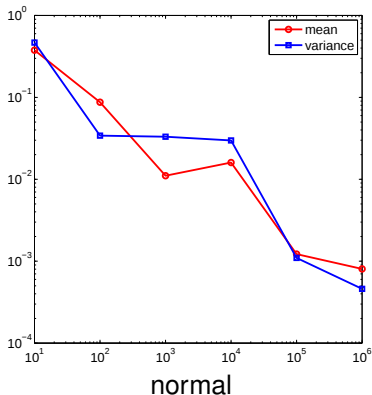
- the Monte Carlo integration error $e_N[\zeta]$ by

$$e_N[\zeta] = E_p[\zeta] - E_N[\zeta] \approx \text{Var}_p[\zeta]^{\frac{1}{2}} N^{-1/2} \nu, \quad \nu \sim N(0, 1)$$

- the error $e_N[\zeta]$ is $O(N^{-1/2})$
- with a constant $\sim$ the variance of the integrand ζ
- the estimate is independent of the dimensionality m



the convergence of mean and variance for one-dim. distr.



- the convergence is not steady, not necessarily monotone
- the approximation of high-order moments is generally harder ...



i.i.d. samples from an implicit and high-dim. distr. are nontrivial.

■ importance sampling

$q(f)$ is an easy-to-sample p.d.f. and close to the posteriori $p(f)$
approximate the expectation of the function ζ w.r.t. $p(f)$ by

$$\int \zeta(f) p(f) df = \int \zeta(f) \frac{p(f)}{q(f)} q(f) df \approx \frac{1}{N} \sum_{i=1}^N \zeta(f^{(i)}) w_i,$$

i.i.d. samples $\{f^{(i)}\}_{i=1}^N$ are drawn from $q(f)$, with $w_i = \frac{p(f^{(i)})}{q(f^{(i)})}$.

The efficiency relies on the approximation quality of $q(f)$ to $p(f)$



Markov chain Monte Carlo (MCMC) general-purposed approach for exploring posteriori $p(f)$

- basic idea: given $p(f)$, construct an **aperiodic and irreducible** Markov chain such that its stationary distribution is $p(f)$.
- By running the chain for **sufficiently long**, simulated values from the chain are **dependent** samples from $p(f)$, and used for computing summarizing statistics.

N. Metropolis, A. Rosenbluth, et al. J. Chem. Phys. 21 (6): 1087–1092, 1953.

W. K. Hastings. Biometrika 57 (1): 97–109, 1970.



The Metropolis-Hastings algorithm is the most basic MCMC method

```
1: Initialize  $f^{(0)}$  and set  $N$ ;  
2: for  $i = 0 : N$  do  
3:   sample  $u \sim U(0, 1)$ ;  
4:   sample  $f^{(*)} \sim q(f^{(i)}, f^{(*)})$   
5:   if  $u < \alpha(f^{(i)}, f^{(*)})$  then  
6:      $f^{(i+1)} = f^{(*)}$ ;  
7:   else  
8:      $f^{(i+1)} = f^{(i)}$ ;  
9:   end if  
10: end for
```

- the uniform distribution $U(0, 1)$
- $p(f)$: the target distribution
- $q(f, f')$ is an easy-to-sample proposal distribution



Given the candidate $f' \sim q(f, f')$, accept it as the new state of the chain with probability $\alpha(f, f')$:

$$\alpha(f, f') = \min \left\{ 1, \frac{p(f')q(f', f)}{p(f)q(f, f')} \right\}.$$

if we reject f' , then the chain remains in the current state f .

- $p(f)$ enters only through α via the ratio $p(f')/p(f)$
 $\Rightarrow$ normalizing constant not needed
- if q is symmetric, i.e., $q(f, f') = q(f', f)$, $\alpha(f, f')$ reduces to

$$\alpha(f, f') = \min \left\{ 1, \frac{p(f')}{p(f)} \right\}.$$

The chain converges to $p(f)$ for **any reasonable $q(f, f')$**



random walker sampler (chain driven by random walk)

- If $q(f, f') = \theta(f' - f)$ for p.d.f. θ , then $f^{(*)} = f^{(i)} + \xi$, $\xi \sim \theta$
- θ : uniform, multivariate normal or t -distribution
- With i.i.d. Gaussian distribution $N(0, \sigma^2)$,

$$f_j^{(*)} = f_j^{(i)} + \xi, \quad \text{with } \xi \sim N(0, \sigma^2)$$

σ^2 controls the size of the random walks, and should be carefully tuned to improve the MCMC convergence.

Heuristic: the optimal acceptance ratio ~ 0.25



independent sampler $q(f, f') = q(f')$

- the acceptance probability $\alpha(f, f')$

$$\alpha(f, f') = \min\{1, w(f')/w(f)\},$$

$w(f) = p(f)/q(f)$ is the importance weight function.

- how to generate $q(f)$: Gaussian approx. by linearized forward model, coarse-scale/reduced-order representation

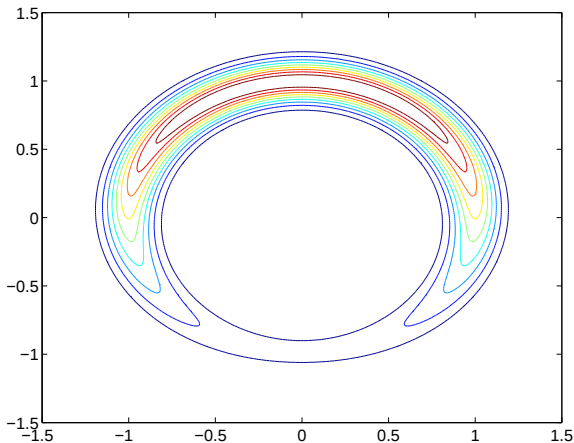
implementation issues

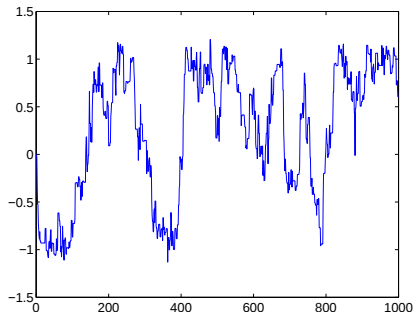
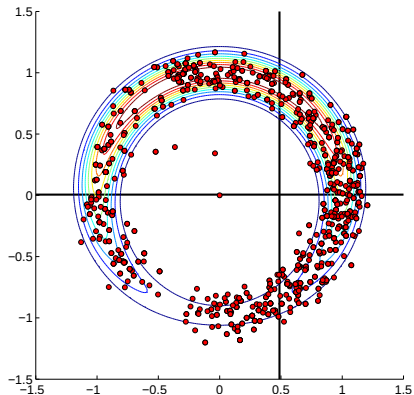
- the first samples are poor approximations as samples from $p(f)$
- discards these initial samples (burning-in)
- assess the convergence of the MCMC chains



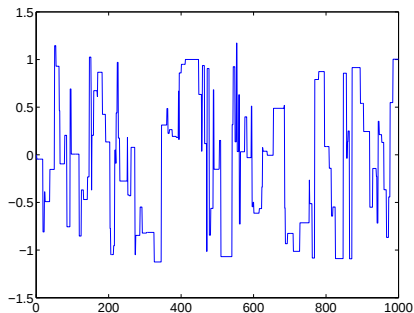
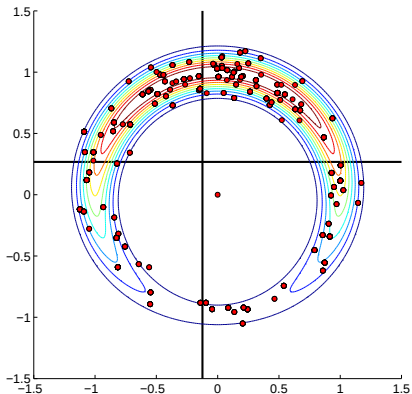
sample the posteriori

$$p(f) \propto e^{-50((f_1^2 + f_2^2)^{1/2} - 1)^2 - 0.5(f_2 - 1)^2}$$

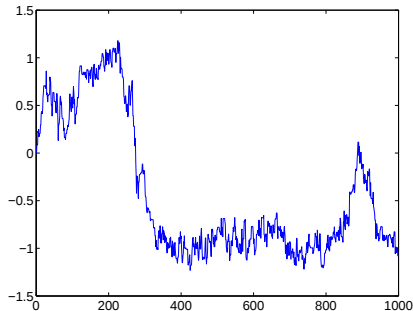
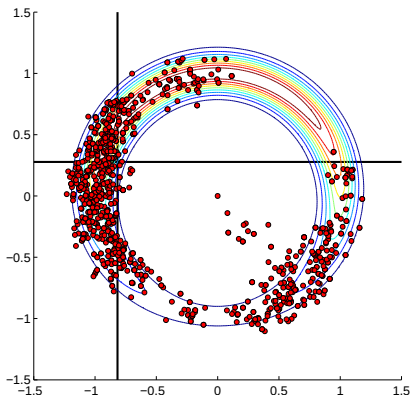




$\sigma = 0.2$, 1000 samples

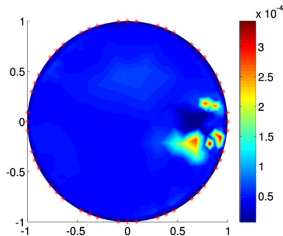
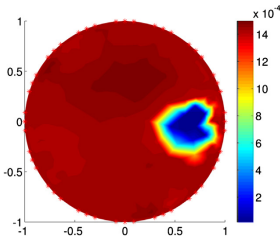
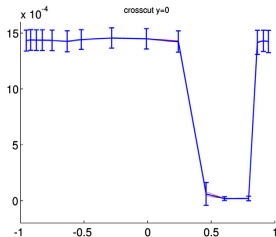
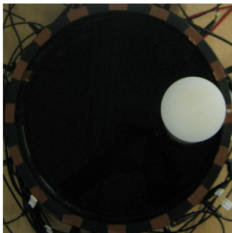


$$\sigma = 1$$



$$\sigma = 0.1$$

■ the proper choice of σ is very crucial ...



230 unknowns, 1×10^7 samples, approximately one day



- If the state space is high dim., it is difficult to update the entire vector f in one single step since $\alpha(f, f')$ is often very small.
- to update only a part of f each time and to implement an updating cycle inside each step
- Gibbs sampler: update one component each time

to update f_i of f , proposal $q(f, f')$: the full conditional

$$q(f, f') = \begin{cases} p(f'_i | f_{-i}) & f'_i = f_{-i}, \\ 0 & \text{otherwise,} \end{cases}$$

where f_{-i} denotes $(f_1, \dots, f_{i-1}, f_{i+1}, \dots, f_m)^t$.

S. Geman, D. Geman. IEEE Trans. Pattern Anal. Mach. Int., 11(6), 721-741, 1984



With this proposal, the acceptance probability $\alpha(f, f')$ is given by

$$\begin{aligned}\alpha(f, f') &= \frac{p(f')q(f', f)}{p(f)q(f, f')} = \frac{p(f')/p(f'_i|f_{-i})}{p(f)/p(f_i|f'_{-i})} \\ &= \frac{p(f')/p(f'_i|f'_{-i})}{p(f)/p(f_i|f_{-i})} = \frac{p(f'_{-i})}{p(f_{-i})} = 1,\end{aligned}$$

these proposals are automatically accepted.



Gibbs algorithm

- 1: Initialize $f^{(0)}$ and set N .
- 2: **for** $i = 0 : N$ **do**
- 3: sample $f_1^{(i+1)} \sim p(f_1 | f_2^{(i)}, f_3^{(i)}, \dots, f_m^{(i)})$,
- 4: sample $f_2^{(i+1)} \sim p(f_2 | f_1^{(i+1)}, f_3^{(i)}, \dots, f_m^{(i)})$,
- 5: $\vdots$
- 6: sample $f_m^{(i+1)} \sim p(f_m | f_1^{(i+1)}, f_2^{(i+1)}, \dots, f_{m-1}^{(i+1)})$,
- 7: **end for**



example: Gibbs sampler for Gaussian noise + smoothness prior
 $p(\lambda) \propto \lambda^{a_0-1} e^{-b_0\lambda}$ on the scale parameter λ , i.e., posteriori

$$p(f, \lambda) \propto e^{-\frac{\tau}{2} \|Af - g\|^2} \cdot \lambda^{\frac{m}{2}} e^{-\frac{\lambda}{2} f^t W f} \lambda^{a_0-1} e^{-b_0\lambda},$$

W encodes the local interaction structure

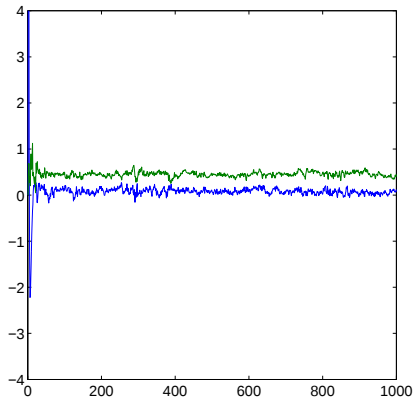
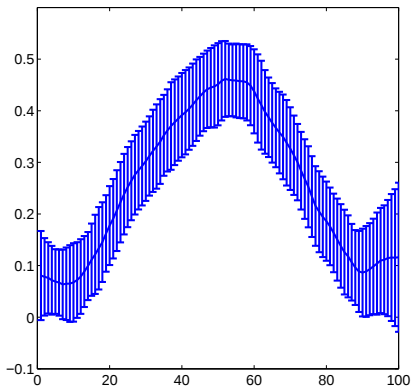
$$p(f_i | f_{-i}, \lambda) \sim N(\mu_i, \sigma_i^2), \quad \mu_i = \frac{b_i}{2a_i}, \quad \sigma_i = \frac{1}{\sqrt{a_i}},$$

with a_i and b_i given by

$$a_i = \tau \sum_{j=1}^n A_{ji}^2 + \lambda W_{ii} \quad \text{and} \quad b_i = 2\tau \sum_{j=1}^n \mu_j A_{ji} - \lambda \mu_p,$$

and $\mu_j = g_j - \sum_{k \neq i} A_{jk} f_k$ and $\mu_p = \sum_{j \neq i} W_{ji} f_j + \sum_{k \neq i} W_{ik} f_k$
 full conditional for λ :

$$p(\lambda | f) \sim G\left(\lambda; \frac{m}{2} + a_0, \frac{1}{2} f^t W f + \beta_0\right).$$



100-dimensional problem, 1000 samples, a few seconds



Approximate inference

Markov chain Monte Carlo: universal technique, asymptotically exact, but generally expensive, convergence diagnosis

approximate inference methods

- variational Bayes Beal 2003; Jordan et al 1999
- expectation propagation Minka 2001...
- Stein variational gradient descent Liu-Wang 2016
- Laplace approximation
- Monte Carlo drop out
-

much faster, but can be limited in accuracy ...



variational Bayes ...

- main idea: to transform the inference problem of exploring $p(f|g)$ into an equivalent optimization problem, and solves it approx.
- probabilistic metrics: Kullback-Leibler divergence, L^1 metric, ...
- given two pdfs q and $\tilde{q}$, the Kullback-Leibler divergence $D_{KL}(q, \tilde{q})$ is defined by

$$D_{KL}(q, \tilde{q}) = \int q(f) \ln \frac{q(f)}{\tilde{q}(f)} df.$$

S. Kullback, R. Leibler. Ann. Math. Statistics 22, 79–86, 1951



- $D_{KL}(q, \tilde{q})$ is asymmetric in q and $\tilde{q}$, and does not satisfy the triangle inequality.
- by Jensen's inequality for the convex function $\varphi(f) = -\ln f$, the divergence $D_{KL}(q, \tilde{q})$ is always nonnegative:

$$\begin{aligned} D_{KL}(q, \tilde{q}) &= - \int q(f) \ln \frac{\tilde{q}(f)}{q(f)} df \geq - \ln \int q(f) \cdot \frac{\tilde{q}(f)}{q(f)} df \\ &= - \ln \int \tilde{q}(f) df = - \ln 1 = 0, \end{aligned}$$

and it vanishes if and only if $q = \tilde{q}$ almost everywhere



Example: $q(f) = N(f; \mu, C)$ and $\tilde{q}(f) = N(f; \tilde{\mu}, \tilde{C})$, with $\mu, \tilde{\mu} \in \mathbb{R}^m$ and $C, \tilde{C} \in \mathbb{R}^{m \times m}$ being positive definite, is given by

$$D_{KL}(q, \tilde{q}) = \frac{1}{2} \left[(\mu - \tilde{\mu})^t \tilde{C}^{-1} (\mu - \tilde{\mu}) + \text{tr}(C \tilde{C}^{-1}) - m - \ln \frac{|C|}{|\tilde{C}|} \right],$$

if $C = \tilde{C}$, the Kullback-Leibler divergence $D_{KL}(q, \tilde{q})$ simplifies to

$$D_{KL}(q, \tilde{q}) = \frac{1}{2} (\mu - \tilde{\mu})^t C^{-1} (\mu - \tilde{\mu})$$



- For any pdf. $q(f)$, the normalizing constant $p(g)$ satisfies

$$\ln p(g) \geq \int q(f) \ln \frac{p(f, g)}{q(f)} df.$$

- The lower bound

$$\text{LB}(g, q) := \int q(f) \ln \frac{p(f, g)}{q(f)} df$$

- $D_{KL}(q(f), p(f|g))$: the gap between $\ln p(g)$ and lower bound
 $\max \text{LB} \equiv \min D_{KL}(q(f), p(f|g))$ w.r.t. q .



Why the KL minimization is hard ... coupling ...

- the separability assumption *decouples* the dependence between different factors
⇒ make the approximation computable ⇒ mean field approx
- separability assumption yields explicit solutions for each component in terms of the others
⇒ alternating direction iterative scheme

M.I. Jordan, Z. Ghahramani, T.S. Jaakkola, L.K. Saul. Mach. Learn. 37, 183–233, 1999.



the assumption $q(f) = \prod q_i(f_i)$ with f_i being disjoint $\Rightarrow$

$$\begin{aligned} \text{LB} &= \int \prod_i q_i(f_i) \left\{ \ln p(f, g) - \sum_i \ln q_i(f_i) \right\} df \\ &= \int q_j(f_j) \left\{ \int \ln p(f, g) \prod_{i \neq j} q_i(f_i) df_{-j} \right\} df_j \\ &\quad - \int q_j(f_j) \ln q_j(f_j) df_j + T(f_{-j}), \end{aligned}$$

define a joint density $\tilde{p}(f_j, g) \propto e^{E_{q_{-j}(f_{-j})}[\ln p(f, g)]}$, $\Rightarrow$

$$\text{LB} = \int q_j(f_j) \ln \frac{\tilde{p}(f_j, g)}{q_j(f_j)} df_j + T(f_{-j}).$$

the optimal $q_j(f_j)$ is then given by

$$q_j(f_j) = \tilde{p}(f_j|g) \propto e^{E_{q_{-j}(f_{-j})}[\ln p(f, g)]},$$



Variational Bayes under the product density restriction

- 1: Initialize $q_i(f_i), i = 2, \dots, M$.
- 2: **for** $k = 1 : K$ **do**
- 3: Update sequentially $q_j(f_j)$ by

$$q_j^{k+1}(f_j) = \frac{e^{E_{q_{-j}^k(f_{-j})}[\ln p(f, g)]}}{\int e^{E_{q_{-j}^k(f_{-j})}[\ln p(f, g)]}}$$

- 4: Check the stopping criterion.
- 5: **end for**

- the convexity of KLD $\Rightarrow$ the convergence to a local optimum
- especially suited to conjugate families



Example: linear inversion + sparsity constraints

- scale mixture representation for the Laplace prior

$$\frac{\lambda}{2} e^{-\lambda|z|} = \int_0^\infty \frac{1}{\sqrt{2\pi s}} e^{-\frac{z^2}{2s}} \cdot \frac{\lambda^2}{2} e^{-\frac{\lambda^2}{2}s} ds.$$

⇒ hierarchical representation of the Bayesian model

$$p(f, w|g) \propto e^{-\frac{\tau}{2} \|Af-g\|^2} \prod_i w_i^{-\frac{1}{2}} e^{-\frac{f_i^2}{2w_i}} \frac{\lambda^2}{2} e^{-\frac{\lambda^2}{2} w_i}.$$

- obstruction: the strong coupling between x and w .
- mean-field approximation on x and w ⇒

$$\begin{aligned} \ln q^k(f) &\propto E_{q^{k-1}(w)}[\ln p(f, w, g)], \\ \ln q^k(w) &\propto E_{q^k(f)}[\ln p(f, w, g)]. \end{aligned}$$

■ Gaussian distribution

$$\begin{aligned} E_{q^{k-1}(w)}[\ln p(f, \lambda, g)] &= E_{q^{k-1}(w)}\left[-\frac{\tau}{2}\|Af - g\|^2 - \sum \frac{1}{2w_i} f_i^2\right] + T(w) \\ &= -\frac{\tau}{2}\|Af - g\|^2 - \frac{1}{2} f^t W f + T(w), \end{aligned}$$

$$\text{with } [W]_{ii} = E_{q^{k-1}(w_i)}\left[\frac{1}{w_i}\right]$$

■ generalized inverse Gaussian

$$\begin{aligned} E_{q^k(f)}[\ln p(f, w, g)] &= \sum_i E_{q^k(f)}\left[-\frac{1}{2} \ln w_i - \frac{f_i^2}{2w_i} - \frac{\lambda^2}{2} w_i\right] + T(f) \\ &= \sum_i \left[-\frac{1}{2} \ln w_i - \frac{1}{2w_i} E_{q^k(f)}[f_i^2] - \frac{\lambda^2}{2} w_i\right] + T(f), \end{aligned}$$

$$q^k(w) = \prod_i \text{GIG}(w_i; \frac{1}{2}, \lambda^2, E_{q^k(f)}[f_i^2]).$$

these two updates that the iterates $q^k(f)$ and $q^k(w)$ would remain Gaussian and (generalized) inverse Gaussian, respectively



outlook

- different perspectives on linear inversion
- there are still many unknown/unclear situations computationally / theoretically
- theoretical understanding of approximate inference techniques ...