

M.PHIL. IN STATISTICAL SCIENCE

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Wednesday 4 June 2008 9.00 to 12.00

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**TIME SERIES AND MONTE CARLO INFERENCE**

*Attempt no more than **FOUR** questions.*

*There are **SIX** questions in total.*

*The questions carry equal weight.*

**STATIONERY REQUIREMENTS    SPECIAL REQUIREMENTS**

*Cover sheet*

*None*

*Treasury tag*

*Script paper*

<p><b>You may not start to read the questions printed on the subsequent pages until instructed to do so by the Invigilator.</b></p>
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## 1 Time Series

Explain what is meant by a *weakly stationary process*  $\{X_t\}$ . Define the *autocovariance function* and the *autocorrelation function* of  $\{X_t\}$ .

Let

$$X_t = \alpha(X_{t-1} - X_{t-2}) + \epsilon_t, \quad (1)$$

where  $\alpha$  is a real constant and  $\{\epsilon_t\}$  is a white noise process with mean zero and variance  $\sigma^2$ . Determine the range of possible values of  $\alpha$  for which (1) has a unique weakly stationary solution.

For  $\alpha = -1/12$ , find the Wold representation of  $\{X_t\}$  and determine the autocovariance function of  $\{X_t\}$ .

[Results from lectures may be quoted and used without proof.]

## 2 Time Series

Let  $\{X_t\}_{t \in \mathbb{Z}}$  be a weakly stationary process with autocovariance function  $\gamma_k$  and spectral density function  $f_X(\lambda)$ . Write down an expression for  $\gamma_k$  in terms of  $f_X(\lambda)$ .

The process  $\{Y_t\}$  is obtained from  $\{X_t\}$  by applying the filter  $\{a_r\}_{r \in \mathbb{Z}}$ , with  $a_r \in \mathbb{R}$  for all  $r \in \mathbb{Z}$  and  $\sum_{r \in \mathbb{Z}} |a_r| < \infty$ , so that  $Y_t = \sum_{r \in \mathbb{Z}} a_r X_{t-r}$ . Show that  $\{Y_t\}$  is weakly stationary and find its spectral density function  $f_Y(\lambda)$  in terms of  $f_X(\lambda)$  and  $a(\lambda) = \sum_{r \in \mathbb{Z}} a_r e^{ir\lambda}$ .

Let  $\{Z_t\}$  be obtained from  $\{Y_t\}$  by applying the filter  $\{b_r\}$ , with  $b_r \in \mathbb{R}$  for all  $r \in \mathbb{Z}$  and  $\sum_{r \in \mathbb{Z}} |b_r| < \infty$ . Write down the spectral density function  $f_Z(\lambda)$  of  $\{Z_t\}$ . Show that  $\{Z_t\}$  can be obtained from  $\{X_t\}$  by applying a linear filter  $\{c_r\}$ , and find  $c_r$  in terms of the  $a_k$ 's and the  $b_k$ 's.

Let the *gain* of a filter  $\{a_r\}$  be  $G_a(\lambda) = |a(\lambda)|$ ,  $\lambda \in [0, \pi]$ .

- (a) Suppose that  $Y_t = X_t - X_{t-1}$ . Find  $f_Y(\lambda)$ . Sketch the gain of the filter taking  $\{X_t\}$  to  $\{Y_t\}$  and comment.
- (b) Suppose that  $Z_t = Y_t - Y_{t-12}$ . Find  $f_Z(\lambda)$ . Sketch the gain of the filter taking  $\{Y_t\}$  to  $\{Z_t\}$  and comment.
- (c) Find the filter that takes  $\{X_t\}$  onto  $\{Z_t\}$  and find its gain.

**3 Monte Carlo Inference**

- (a) Monte Carlo methods depend crucially on the ability to generate pseudo random numbers in the interval  $(0, 1)$ . Make a short list of what you consider to be the most important properties of a good pseudo random number generator.
- (b) Suppose you had an ideal pseudo random number generator giving you the ability to generate arbitrarily many independent uniform variates  $U_1, U_2, \dots$ , i.e.,  $U_i \sim U(0, 1)$ .
- (i) Describe how the method of inversion can be used to obtain draws from the  $\text{Exp}(\lambda)$  distribution. How can this method be extended to obtain draws from a double-exponential (Laplace) distribution? What property of the pseudo random number generator is crucial to ensure that your algorithm gives samples from the correct distribution?
  - (ii) Give two distinct algorithms for obtaining draws from a  $\chi_\nu^2$  distribution where  $\nu \in \{2, 3, \dots\}$ . Say which algorithm you prefer, and why. Your answer may depend on  $\nu$ .
  - (iii) Consider obtaining draws from a  $\text{Beta}(\alpha, \beta)$  distribution, when  $\alpha, \beta \in \{1, 2, \dots\}$ . Give one algorithm which uses one of the methods from (ii), and one new method.

#### 4 Monte Carlo Inference

- (a) Describe the jackknife *and* nonparametric bootstrap methods for estimating the variance of an estimator  $\hat{\theta}$  of some parameter  $\theta(F)$ , on the basis of a random sample  $x_1, \dots, x_n$  of distinct observations from  $F$ . Your description of the nonparametric bootstrap should include the form of the empirical distribution function  $\hat{F}_n$  used in the algorithm.
- (b) Find the probability that a bootstrap sample contains at least one repeated value.
- (c) Consider the following R code where  $\mathbf{x}$  is a vector of length  $n$  containing the random sample  $x_1, \dots, x_n$ , where  $\mathbf{x}$  and  $n$  have been set earlier in the code.

```
R1a> mat <- matrix(NA, nrow=n, ncol=n-1)
R2a> for(i in 1:n) mat[i,] <- x[-i]
R3a> vect <- apply(mat, 1, mean)
R4a> (n-1)*mean((vect - mean(vect))^2)
R5a> (n-1)*(mean(vect) - mean(x))
```

Explain what is being calculated in lines R4a and R5a. Give the numerical value of the expression in line R5a, and justify your answer.

Now consider another piece of R code below (with the same  $\mathbf{x}$  as above).

```
R1b> alpha <- 0.05
R2b> B <- 199
R3b> mat <- matrix(NA, nrow=B, ncol=n)
R4b> for(b in 1:B) mat <- sample(x, n, replace=TRUE)
R5b> vect <- apply(mat, 1, mean)
R6b> s <- sort(vect)
R7b> c(s[(B+1)*alpha/2], s[(B+1)*(1-alpha/2)])
```

Explain what is being calculated in the code, with particular attention paid to the value of the expression in line R7b.

- (d) Suppose that we had another random sample  $y_1, \dots, y_m$  from a distribution  $G \neq F$ , where  $\theta = \mathbb{E}_F\{X\} = \mathbb{E}_G\{Y\}$  and  $\text{Cov}(X, Y) < 0$ . Give an algorithm for constructing an efficient, unbiased estimator  $\tilde{\theta}$  of  $\theta$  that uses the combined sample  $x_1, \dots, x_n, y_1, \dots, y_m$ . How could you estimate  $\text{Var}(\tilde{\theta})$ ?

## 5 Monte Carlo Inference

- (a) (i) Describe the Gibbs Sampler for obtaining a dependent sample from some distribution  $\pi(\boldsymbol{\theta})$ ,  $\boldsymbol{\theta} \in \mathbb{R}^p$ .
- (ii) Suppose that we observe data  $\mathbf{y} = (y_1, \dots, y_n)^\top$ , with corresponding known (scalar) covariates  $\mathbf{x} = (x_1, \dots, x_n)^\top$  and that we want to fit a polynomial regression model of order  $k$  to the data. Then we can express the model in the form

$$\mathbf{y} = \mathbf{X}_k \boldsymbol{\beta}_k + \boldsymbol{\varepsilon}$$

for design matrix

$$\mathbf{X}_k = \begin{pmatrix} 1 & x_1 & \cdots & x_1^k \\ \vdots & \vdots & & \vdots \\ 1 & x_n & \cdots & x_n^k \end{pmatrix}$$

where  $\boldsymbol{\beta}_k = (\beta_0, \beta_1, \dots, \beta_k)^\top$  and  $\boldsymbol{\varepsilon} = (\varepsilon_1, \dots, \varepsilon_n)^\top$ , with  $\boldsymbol{\varepsilon} \sim N_n(\mathbf{0}, \sigma^2 \mathbf{I})$ , where  $\mathbf{I}$  is the  $n \times n$  identity matrix. For independent priors  $\sigma^2 \sim \Gamma^{-1}(a, b)$  and  $\boldsymbol{\beta}_k \sim N_{k+1}(\boldsymbol{\mu}_k, \boldsymbol{\Sigma}_k)$  the posterior distribution is given by

$$\begin{aligned} \pi(\boldsymbol{\beta}_k, \sigma^2 | \mathbf{x}, \mathbf{y}) &\propto (\sigma^2)^{-n/2} \exp \left\{ -\frac{1}{2\sigma^2} (\mathbf{y} - \mathbf{X}_k \boldsymbol{\beta}_k)^\top (\mathbf{y} - \mathbf{X}_k \boldsymbol{\beta}_k) \right\} \\ &\times (\sigma^2)^{-(a+1)} \exp \left\{ -\frac{b}{\sigma^2} \right\} \times \exp \left\{ -\frac{1}{2} (\boldsymbol{\beta}_k - \boldsymbol{\mu}_k)^\top \boldsymbol{\Sigma}_k^{-1} (\boldsymbol{\beta}_k - \boldsymbol{\mu}_k) \right\}. \end{aligned}$$

Show that the conditional distributions  $\pi(\boldsymbol{\beta}_k | \sigma^2, \mathbf{x}, \mathbf{y})$  and  $\pi(\sigma^2 | \boldsymbol{\beta}_k, \mathbf{x}, \mathbf{y})$  are multivariate normal and inverse gamma, respectively, and calculate the parameters of each distribution.

- (iii) Hence describe how we can use the Gibbs Sampler to obtain a dependent sample from the joint posterior distribution of  $\pi(\boldsymbol{\beta}, \sigma^2 | \mathbf{x}, \mathbf{y})$ .
- (b) Now suppose that the order of the polynomial is unknown, and that we wish to use a reversible jump procedure to update the order of the polynomial model. We propose to move from the model of order  $k$ , with parameters  $\boldsymbol{\beta}_k$ , to the model of order  $k+1$  with parameters  $\boldsymbol{\beta}'_{k+1}$  (keeping  $\sigma^2$  fixed) using the following procedure,

$$\begin{aligned} \beta'_i &= \beta_i && \text{for } i = 1, \dots, k \\ \beta'_{k+1} &= z && \text{for } z \sim N(0, \sigma_\beta^2) \text{ and } \sigma_\beta^2 \text{ known} \\ \beta'_0 &= \beta_0 - \frac{z}{n} \sum_{i=1}^n x_i^{k+1}. \end{aligned}$$

- (i) Calculate an explicit expression for the corresponding acceptance probability for this move.
- (ii) Define the reverse move, for moving from the model of order  $k+1$  to the model of order  $k$ .
- (iii) What is the corresponding acceptance probability for this reverse move, from the model of order  $k+1$ , to the model of order  $k$ ?

## 6 Monte Carlo Inference

- (a) Let  $\mathbf{x}$  represent observed data, and  $\mathbf{z}$  denote missing data, with joint distribution  $f(\mathbf{x}, \mathbf{z}; \boldsymbol{\theta})$ . Briefly describe the iterative Expectation Maximisation (EM) algorithm for finding the  $\hat{\boldsymbol{\theta}}$  that maximises the observed data likelihood  $L(\mathbf{x}|\boldsymbol{\theta})$ .
- (b) Suppose that  $\mathbf{y} = (y_1, y_2, y_3, y_4)$  is a data vector of observed counts from a multinomial distribution with parameters  $n$  and  $\mathbf{p}$ , where the cell probabilities

$$\mathbf{p} = (p_1, p_2, p_3, p_4) = \left( \frac{1}{2} - \frac{\theta}{2}, \frac{\theta}{4}, \frac{\theta}{4}, \frac{1}{2} \right),$$

are parameterised by  $\theta \in [0, 1]$ .

- (i) Find the maximum likelihood estimator  $\hat{\theta}$  based on the complete data likelihood  $L(\theta|\mathbf{y})$ .
- (ii) Now suppose instead that only three counts

$$\mathbf{x} = (x_1, x_2, x_3)$$

were observed, where

$$\mathbf{x} = (y_1, y_2, y_3 + y_4).$$

That is,  $y_3 = x_3 - z$ ,  $y_4 = z$  and  $z$  is missing. Consider using the EM algorithm for estimating  $\hat{\theta}$  based on the observed data log likelihood  $L(\mathbf{x}|\theta)$ . Derive the “E-step” of the EM algorithm and write the resulting expression(s) in terms of  $\log L(\mathbf{y}^{(t)}|\theta)$ , for

$$\mathbf{y}^{(t)} = (y_1, y_2, y_3^{(t)}, y_4^{(t)})$$

where  $y_3^{(t)} = x_3 - z^{(t)}$ ,  $y_4^{(t)} = z^{(t)}$ , and  $z^{(t)} = \mathbb{E}\{z|\mathbf{x}, \theta^{(t)}\}$  which you should calculate. In other words, show that the “E-step” is the same as “filling in the missing values” in this case.

- (iii) Combine the “E-step” in part (ii) with an “M-step” derived from the appropriate application of your result from part (i). That is, give a complete description of your EM algorithm in this case for iteratively finding  $\tilde{\theta}$ , the maximum likelihood estimator of the observed data likelihood.
- (iv) Suppose that  $\mathbf{x} = (38, 34, 125)$  and  $\theta^{(t)} = 0.5$ . What are the values of  $y_3^{(t)}$  and  $\theta^{(t+1)}$ ?

**END OF PAPER**