SMOOTHNESS PRIORS IN TIME SERIES

BY

WILL GERSCH and GENSHIRO KITAGAWA

TECHNICAL REPORT NO. 391
JUNE 2, 1987

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1. INTRODUCTION ............................................................................................................. 1
2. SMOOTHNESS PRIORS MODELING, LEAST SQUARES ANALYSIS ......................... 5
  2.1 SMOOTHNESS PRIORS BAYESIAN MODELING .................................................. 5
  2.2 SMOOTHNESS PRIORS IN THE FREQUENCY DOMAIN .................................... 9
  2.3 DISCUSSION ......................................................................................................... 14
3. STATE SPACE GAUSSIAN SMOOTHNESS PRIORS MODELING .......................... 16
  3.1 NONSTATIONARY MEAN SMOOTHNESS PRIORS STATE SPACE MODELING ........................................................................................................... 16
  3.2 RECURSIVE ESTIMATION OF STATE AND LIKELIHOOD COMPUTATION .... 19
  3.3 NONSTATIONARY COVARIANCE MODELING .................................................. 21
  3.4 EXAMPLES .......................................................................................................... 25
  3.5 COMMENTS, DISCUSSION .................................................................................. 28
4. STATE SPACE NON GAUSSIAN MODELING OF NONSTATIONARY TIME SERIES .................................................................................................. 31
  4.1 THE NON GAUSSIAN STATE SPACE MODEL .................................................. 32
  4.2 NUMERICAL EXAMPLES ................................................................................... 34
  4.3 COMMENTS, DISCUSSION .................................................................................. 36
5. SUMMARY ................................................................................................................. 38
REFERENCES .................................................................................................................. 39
1. INTRODUCTION

Several different kinds of stationary and nonstationary time series modeling problems are considered here from a Bayesian-smoothness priors approach. The smoothness priors specify the prior distribution of the time series model parameters.

The term "smoothness priors" is very likely due to Shiller (1973). Shiller modeled the distributed lag (impulse response) relationship between the input and output of economic time series under difference equation "smoothness" constraints on the distributed lags. A tradeoff of the goodness-of-fit of the solution to the data and the goodness-of-fit of the solution to a smoothness constraint was determined by a single smoothness tradeoff parameter. Shiller did not offer an objective method of choosing the smoothness tradeoff parameter. Akaike, (1980), completed the analysis initiated by Shiller. Akaike developed and exploited the concept of the likelihood of the Bayesian model and used a maximization of the likelihood procedure for determining the smoothness tradeoff parameter. (In Bayesian terminology, the smoothness tradeoff parameter is referred to as the "hyperparameter", Lindley and Smith, 1972.) The origin of Shiller-Akaike smoothness priors can be seen in a smoothing problem posed by Whittaker (1923). The smoothing problem context is now understood to be common to a variety of other statistical data analysis problems including density estimation and image analysis (Titterington 1985).

In the problem treated by Whittaker, the observations \( y_n, n = 1, \ldots, N \) are given. They are assumed to consist of the sum of a "smooth" function and observation noise,

\[
y_n = f_n + \epsilon_n. \tag{1.1}
\]

The problem is to estimate the unknown \( f_n, n = 1, \ldots, N \). In a time series interpretation of this problem, \( f_n, n = 1, \ldots, N \) is the trend of a nonstationary mean time series. A typical approach to this problem is to use a class of parametric models. The quality of the analysis is completely dependent upon the appropriateness of the assumed model class. A flexible model is desirable. In this context, Whittaker suggested that the solution balance a tradeoff of goodness of fit to the data and goodness of fit to a smoothness criterion. This idea was realized by minimizing
\[
\sum_{n=1}^{N} (y_n - f_n)^2 + \mu^2 \sum_{n=1}^{N} (\nabla^k f_n)^2
\] (1.2)

for some appropriately chosen smoothness tradeoff parameter \(\mu^2\). In (1.2) \(\nabla^k f_n\) expresses a kth-order difference constraint on the solution \(f\), with \(\nabla f_n = f_n - f_{n-1}\), \(\nabla^2 f_n = \nabla(\nabla f_n)\), etc. (Whittaker's original solution was not expressed in a Bayesian context. Whittaker and Robinson (1924) does invoke a Bayesian interpretation of this problem.)

The properties of the solution to the problem in (1.1)-(1.2) are clear. If \(\mu^2 = 0\), \(f_n = y_n\) and the solution is a replica of the observations. As \(\mu^2\) becomes increasingly large, the smoothness constraint dominates the solution and the solution satisfies a kth order constraint. For large \(\mu^2\) and \(k=1\), the solution is a constant, for \(k=2\), it is a straight line etc. Whittaker left the choice of \(\mu^2\) to the investigator.

Kohn and Ansley (1987) demonstrate that the signal extraction problem of (1.1) and the smoothing problem of (1.2) are equivalent problem statements. The equivalence also holds for broad variations of signal extraction and smoothing problems. All of the time series analysis problems that we treat here are variations of the signal extraction/smoothing problem in (1.1) and (1.2).

An implication of Akaike (1980) is that a Bayesian interpretation of the smoothing problem in (1.2) implies that the difference equation constraint is a stochastically perturbed zero-mean unknown variance difference equation. The stochastically perturbed difference equation constraint in the trend estimation problem is a smoothness priors constraint in the time domain. Akaike (1980), considered other time domain smoothness priors constraint problems including the Shiller distributed lag problem and the seasonal adjustment of time series. Ishiguro et al. (1981) used time domain smoothness priors constraints and fixed effects regression in an analysis of tidal effects. Akaike (1979) employed a frequency domain smoothness priors constraint on the distributed lag parameters in the Shiller problem. Gersh and Kitagawa (1984) and Kitagawa and Gersh (1985a) are other frequency domain smoothness priors time series problem analyses.
Shiller (1973), Akaike (1980), and all of the aforementioned smoothness priors analyses, are Bayesian analyses of the linear model with Gaussian stochastic constraints and Gaussian disturbances. The critical ideas in smoothness priors are the likelihood of the Bayesian model and the use of likelihood as a measure of the goodness of fit of the model. In our analysis, hyperparameters have interpretations as noise to signal ratios and they have a remarkable role in the analysis. The maximization of the likelihood of a small number of hyperparameters permits the robust modeling of a time series with relatively complex structure and a very large number of implicitly inferred parameters. When we consider alternative smoothness priors models, with different distributional assumptions or different numbers of parameters to model the same data, we use the Akaike AIC statistic (Akaike 1973), to choose between candidate models. Kitagawa (1987) is a smoothness priors state space modeling of nonstationary time series in which neither the system noise or the observation noise are necessarily Gaussian distributed.

The original Whittaker problem has also given rise to work on splines in numerical analysis and to related smoothing problem analysis, particularly by Wahba (1977),(1982). Smoothness priors relates to the ill-posed problems and problems of statistical regularization that have been considered extensively in the Soviet Union by Tikhonov (1965) and his associates. Also related are the "bump hunting"-penalized likelihood methods, Good and Gaskins (1980) and Wecker and Ansley (1983) and O'Sullivan et al. (1986). Vigorous work, primarily at the Institute of Statistical Mathematics, Tokyo, has resulted in the application of smoothness priors methods to a variety of applications, other than the ones we discuss here. These applications include the seasonal adjustment of time series, (Akaike 1980b), tidal analysis (Ishiguro et al. 1981), binary regression (Ishiguro and Sakamoto 1983), cohort analysis (Nakamura 1986), and density estimation (Tanabe and Sagae 1987).

Smoothness priors problems that are amenable to analysis by least squares algorithms are treated in Section 2. The likelihood of the Bayesian model, as done by Akaike, is in Section 2.1. The Bayesian solution to the smoothing problem originally posed by Whittaker is also shown
there. In that problem, the smoothness priors constraints are time domain constraints. The priors are expressed as zero-mean unknown variance stochastically perturbed kth-order random walk difference equations. In Section 2.2, the estimation of the power spectral density from short duration stationary time series illustrates the use of frequency domain smoothness priors constraints. In Section 3 several smoothness priors nonstationary time series problems, which are amenable to a Kalman filter state space method of analysis, including examples of the modeling of nonstationary mean and nonstationary covariance time series, are shown. All of the aforementioned treat a linear model, Gaussians distributions situation. That method is generalized in Section 4. There we show a smoothness priors state space not necessarily Gaussian nonstationary time series analysis method. Finally, Section 5 is a summary.
2 SMOOTHNESS PRIORS MODELING: LEAST SQUARES ANALYSIS

In Section(2.1) we review the concept of smoothness priors Bayesian modeling as introduced in Akaike (1980). That method is applied to the problem addressed by Whittaker (1923), the estimation of a trend in white noise. The smoothness priors constraint is expressed as a kth order random walk with a zero-mean, unknown variance perturbation. The variance is a hyperparameter of the prior distribution. The constraint is a time domain constraint on the priors. In Section 2.2 we introduce the notion of frequency domain constraint on the priors. That method is used in the estimation of the power spectral density of a stationary time series. Section 2.3 is a discussion.

The frequency domain smoothness priors method is particularly suited for the situation in which only a short span of data is available for analysis. In that case, the results of conventional parametric model analysis methods are particularly sensitive to the choice of model order. We circumvent that problem, using the frequency domain smoothness priors, by tending to fit models that are "too long". Those priors reflect the integrated squared kth derivative with respect to frequency of the departure from model smoothness. The estimation of the model parameters and an additional small number of hyperparameters is required. The maximization of the likelihood of the hyperparameters is the critical computation.

2.1 SMOOTHNESS PRIORS BAYESIAN MODELING

Consider the linear regression model subject to Bayesian-stochastic constraints

\[ \begin{pmatrix} y \\ \theta \end{pmatrix} \sim N \left( \begin{bmatrix} X \theta \\ 0 \end{bmatrix}, \begin{bmatrix} \sigma^2 I & 0 \\ 0 & \lambda^2 D^{-1} D^{-T} \end{bmatrix} \right) \]  

(2.1.1)

The dimensions of the matrices in (2.1.1) are \( y: n \times 1 \); \( X: n \times p \); \( \theta: p \times 1 \). \( \sigma^2 \) and \( \lambda^2 \) are unknown. \( y \) is the vector of observed data, \( X \) and \( D \) are assumed known. \( \theta \) is the normally distributed prior parameter vector. The observation noise variance is \( \sigma^2 \). In this conjugate family Bayesian situation (Berger 1985), the mean of the posterior normal distribution of the parameter vector \( \theta \) minimizes
\[ \| y - X\theta \|^2 + \lambda^2 \| D\theta \|^2 \] (2.1.2)

If \( \lambda^2 \) were known, the computational problem in (2.1.2) could be solved by an ordinary least squares computation. The solution for \( \hat{\theta} \), the posterior mean, is the minimizer of

\[ \| (y) - (X\lambda D)\theta \|^2 \] (2.1.3)

That solution is

\[ \hat{\theta} = \left( X^TX + \lambda^2 D^TD \right)^{-1} X^Ty \] (2.1.4)

with the residual sum of squares,

\[ SSE(\hat{\theta}, \lambda^2) = y^Ty - \hat{\theta}^T X^TX + \lambda^2 D^TD \left( \left( X^TX + \lambda^2 D^TD \right)^{-1} \right) \hat{\theta} \] (2.1.5)

For a smoothness priors interpretation of the problem in (2.1.1) and (2.1.2), multiply (2.1.2) by \(-1/2\sigma^2\) and exponentiate. Then the \( \theta \) that minimizes (2.1.2) also maximizes

\[ \exp \left( -\frac{1}{2\sigma^2} \| y - X\theta \|^2 \right) \exp \left( -\frac{\lambda^2}{2\sigma^2} \| D\theta \|^2 \right) \] (2.1.6)

In (2.1.6), the posterior distribution interpretation of the parameter vector \( \theta \) is that it is proportional to the product of the conditional data distribution (likelihood), \( p(y | X, \theta, \sigma^2) \), and a prior distribution, \( \pi(\theta | \lambda^2, \sigma^2) \) on \( \theta \),

\[ \pi(\theta | y, \lambda^2, \sigma^2) \propto p(y | X, \theta, \sigma^2)\pi(\theta | \lambda^2, \sigma^2) \] (2.1.7)

The integration of (2.1.7) yields \( L(\lambda^2, \sigma^2) \), the likelihood for the unknown parameters \( \lambda^2 \) and \( \sigma^2 \),

\[ L(\lambda^2, \sigma^2) = \int_{-\infty}^{\infty} \pi(\theta | y, \lambda^2, \sigma^2) d\theta \] (2.1.8)

I.J. Good (1965) referred to the maximisation of (2.1.8) as a Type II maximum likelihood method. Since \( \pi(\theta | y, \lambda^2, \sigma^2) \) is normally distributed, (2.1.8) can be expressed in the closed form, (Akaike 1980),

\[ \text{6} \]
\[ L(\lambda^2, \sigma^2) = (2\pi \sigma^2)^{-N/2} |\lambda^2 D^T D|^{1/2} \| X^T X + \lambda^2 D^T D \|^{-1/2} \exp \left( \frac{-1}{2\sigma^2} \text{SSE}(\theta, \lambda^2) \right). \] (2.1.9)

The maximum likelihood estimator of \( \sigma^2 \) is

\[ \hat{\sigma}^2 = \frac{\text{SSE}(\hat{\theta}, \lambda^2)}{N}. \] (2.1.10)

It is convenient to work with \(-2\log\) likelihood. Using (2.1.10) in (2.1.9) yields

\[ -2\log L(\lambda^2, \hat{\sigma}^2) = N \log 2\pi + N \log \text{SSE}(\hat{\theta}, \lambda^2)/N + \log \| X^T X + \lambda^2 D^T D \| - \log \| \lambda^2 D^T D \| + N. \] (2.1.11)

A practical way to determine the value of \( \lambda^2 \) for which the \(-2\log\)-likelihood is minimized, is to compute the likelihood for discrete values of \( \lambda^2 \) and search the discrete \(-2\log\) likelihood-hyperparameter space for the minimum. Akaike (1980) is very likely the first practical use of the likelihood of the Bayesian model and the use of the likelihood of the hyperparameters, as a measure of the goodness of fit of a model to data.

Estimating a Trend

Here we return to the original problem posed in the Introduction. We use the notation \( f_n = t_n \)

where \( t_n \) is the trend at time \( n \) to emphasize the fact that we are estimating the mean of a nonstationary mean time series. A critically important observation is that from the stochastic regression or Bayesian point of view, the difference equation constraints in the Whittaker problem are stochastic. That is, \( \nabla t_n = w_n \), with \( w_n \) assumed to be a normally distributed zero-mean sequence with unknown variance \( \tau^2 \). For example for \( k=1 \) and \( k=2 \) those constraints are:

\[ t_n = t_{n-1} + w_n; \] (2.1.12)

\[ t_n = 2t_{n-1} - t_{n-2} + w_n. \]

A parameterization which relates the trend estimation problem to the earlier development in this section is \( \tau^2 = \sigma^2 / \lambda^2 \). Corresponding to the matrix \( D \) in (2.1.2), for \( k=1 \) and \( k=2 \), the smoothness constraints can be expressed in terms of the following \( N \times N \) constraint matrices:

7
In (2.1.13) $\alpha$ and $\beta$ are small numbers that are chosen to satisfy initial conditions.

For fixed $k$ and fixed $\lambda^2$ the least squares solution can be simply expressed in the form of (2.1.13). For example with $k = 2$, the solution $\{t_n, n = 1, \ldots, N\}$ satisfies

\[
\parallel \begin{bmatrix} y \\ 0 \end{bmatrix} - \left( \begin{bmatrix} I \\ \lambda D \end{bmatrix} t \right) \parallel^2
\]

Note that the problem in (2.1.14) is a version of the Bayesian linear stochastic regression problem in (2.1.3) with $\theta = t = (t_1, \ldots, t_N)^T$, $X = I$, the $N \times N$ identity matrix, and $D = D_1$ or $D_2$. From (2.1.3), the solution to (2.1.14), with $D = D_2$, is

\[
t = [I + \lambda^2 D_2^T D_2]^{-1} y
\]

and the value of $SSE(\hat{\theta}, \lambda^2)$ is given by (2.1.5) with $\hat{\theta} = t, X = I, D = D_2$. The smoothing problem expression of (2.2.15) is that the solution vector is $t = (t_{1|N}, \ldots, t_{N|N})^T$. The least squares problem in (2.1.14), with $D = D_2$, is solved for discrete values of $\lambda^2 D$ and the $-2$ log likelihood-hyperparameter space is searched for a minimum. From (2.1.11), the minimized value of $-2$ log likelihood for this problem is:

\[
-2\log L(\hat{\alpha}^2, \hat{\beta}^2) = N\log 2\pi + N\log SSE((t^2, \hat{\lambda}^2)/N) + \log |\hat{\lambda}^2 D_2^T D_2 + I| - \log |\hat{\lambda}^2 D_2^T D_2| + N.
\]

The numerical values of $SSE(\hat{t}^2, \lambda^2)$ and of the determinants in (2.1.16) are transparent in a least squares algorithm analysis of (2.1.14). Since $\lambda = \sigma^2/r$, $\lambda^2$ has a noise-to-signal ratio interpretation. Smaller $\lambda$ corresponds to smoother trends.
An Example of Trend Estimation

We consider the example of an asymmetrically truncated normal density-like function in the presence of additive noise, $N(0, \sigma^2)$. Figure 1a shows the smooth function $t_n, n=1, \ldots, N$ and the superposition of $t_n$ and the additive noise. The problem is: Given the noisy observations $y_n$, estimate the unknown smooth function that is in the noise, i.e. specify $t_n|N, n=1, \ldots, N$. We solved the least squares computational problem in (2.1.14) using the Householder transformation method. -2 log likelihood of the hyperparameter model is computed from (2.1.16).

The critical role of the hyperparameter is transparent in this example. Figures 1b,c,d show the estimated trend for values of the hyperparameters that are too small, ($\lambda^2 = 0.00001$), and too large ($\lambda = 0.1$) as well as the hyperparameter for which -2log likelihood is minimized, ($\lambda = 0.00136$). As anticipated, with the hyperparameter defined as indicated above, the estimated trend for a too large value of the hyperparameter is too bumpy and the estimated trend for a too small value of the hyperparameter is too smooth.

It is important to note that in this example, although the truncated Gaussian satisfies $\nabla^2 \log t_n = 0$, we estimate the trend with the "incorrect" model $\nabla^2 t_n = w_n$, the stochastically perturbed second order difference equation. The point is that a priori we do not know a correct expression for the underlying smooth function in (1.1). Different hyperparameter values result in solutions of the stochastically perturbed second order difference equation with very different smoothness properties. The best of those solutions yields a very good approximation to the original unknown smooth function. This key observation was referred to by Shiller, (1973), as the "flexible ruler approach".

2.2 SMOOTHNESS PRIORS IN THE FREQUENCY DOMAIN

The smoothness priors in the estimation of the mean value of a nonstationary time series was expressed as a time domain, stochastically perturbed difference equation constraint on the evolution of the trend. Smoothness priors contraints can also be expressed in the frequency domain. In
this section, we illustrate the use of frequency domain priors for the estimation of the power spectral density of a stationary time series.

A Long AR Model For Spectral Estimation

A smoothness priors-long autoregressive (AR) model approach is used here for spectral density estimation.

The classical windowed periodogram method of spectral estimation is satisfactory for spectral analysis when the data set is "long". The alternative of spectral estimation via the fitting of parametric models to moderate length data spans became popular in the last decade, Kesler(1986). When the data span is relatively short, three facts render parametric modeling methods of spectral estimation statistically unreliable. One is the instability or small sample variability of whatever statistic is used for determining the best order of parametric model fitted to the data. The second is that usually the "parsimonious" parametric model is not a very good replica of the system that generated the data. The third is that the spectral density of the fitted parametric model can not possibly be correct. Independent of which parametric model order is selected, there is information in the data to select models of different orders. A Bayesian estimate of power spectral density requires that the spectral density of parametric models of different model orders be weighted in accordance with the likelihood and the prior of the model order of different models.

The smoothness priors AR model of spectral estimation alleviates this problem. A particular class of frequency domain smoothness priors is assumed for the coefficients of AR model order M, with M relatively large. The likelihood of the hyperparameters that characterize the class of smoothness priors is maximized to yield the best AR model of order M with the best data dependent priors. (A more complete treatment of the modeling discussed here is in Kitagawa and Gersch, 1985a.)
The Smoothness Priors Long AR Model

Consider the autoregressive model of order $M$,

$$ y_n = \sum_{m=1}^{M} a_m y_{n-m} + \epsilon_n $$  \hspace{1cm} (2.2.1)

In (2.2.1) \{\epsilon_n\} is a Gaussian white noise with mean zero and variance $\sigma^2$. A least squares fit of the AR model to the data, $y_1, \ldots, y_N$, with the first $M$ observations $y_{1-M}, y_{2-M}, \ldots, y_0$ treated as given constants, leads to the minimization of

$$ \sum_{n=1}^{N} |y_n - \sum_{m=1}^{M} a_m y_{n-m}|^2 $$  \hspace{1cm} (2.2.2)

If $M$ is comparable to $N$, the result of the least squares computation can be meaningless. The smoothness priors solution mitigates this difficulty by considering the solution of the constrained least squares problem. We consider a frequency domain smoothness priors constraint on the distribution of the AR model parameters. The frequency response function of the whitening filter of the AR process given by

$$ A(f) = 1 - \sum_{m=1}^{M} a_m \exp \{-2\pi imf\}. $$  \hspace{1cm} (2.2.3)

Let a measure of the smoothness of the frequency response function be

$$ R_k = \int_{-\frac{1}{2}}^{\frac{1}{2}} \left| \frac{d^k A(f)}{df^k} \right|^2 df = (2\pi)^{2k} \sum_{m=1}^{M} m^{2k} a_m^2. $$  \hspace{1cm} (2.2.4)

From the definition in (2.2.4), a large value of $R_k$ means an unsmooth (in the sense of differential) frequency response function. We also use the zero derivative smoothness constraint,

$$ R_0 = \int_{-\frac{1}{2}}^{\frac{1}{2}} |A(f)|^2 df = 1 + \sum_{m=1}^{M} a_m^2 $$  \hspace{1cm} (2.2.5)

as a penalty to the whitening filter.

With these constraints, and with $\lambda^2$ and $\nu^2$ fixed, the AR model coefficients \{a_m, m=1,\ldots,M\}, minimize
\[ \sum_{n=1}^{N} |y_n - \sum_{m=1}^{M} a_m y_{n-m}|^2 + \lambda^2 \sum_{m=1}^{M} m^{2k} a_m^2 + \nu^2 \sum_{m=1}^{M} a_m^2. \]  

(2.2.6)

In (2.2.6), \( \lambda^2 \) and \( \nu^2 \) are the hyperparameters. By a proper choice of these parameters, our estimates of the AR model coefficients balance the tradeoff between the infidelity of the model to the data and the infidelity of the model to the frequency domain smoothness constraints. For completeness, to within a constant, the Gaussian priors on the AR model coefficients corresponding to the \( R_0 \) and \( R_4 \) constraints are

\[
\exp\left(-\frac{\lambda^2}{2} \sigma^2 \sum_{m=1}^{M} m^{2k} a_m^2 \right) \exp\left(-\frac{\nu^2}{2} \sigma^2 \sum_{m=1}^{M} a_m^2 \right).
\]

(2.2.7)

Following our earlier discussion, define the matrices \( D \) and \( a \) and the matrices \( X \) and \( y \) by

\[
D = \begin{bmatrix}
(\nu^2 + \lambda^2)^{1/2} & (\nu^2 + 2M\lambda^2)^{1/2} \\
\vdots & \ddots \\
(\nu^2 + M^{2k}\lambda^2)^{1/2} & a_1 & a_2 & \cdots & a_M \\
& y_0 & y_1 & \cdots & y_N \\
\end{bmatrix}, \quad \quad a = \begin{bmatrix}
a_1 \\
a_2 \\
\vdots \\
a_M \\
\end{bmatrix}, \quad X = \begin{bmatrix}
y_0 & \cdots & y_1 & \cdots & y_M \\
y_1 & \cdots & y_2 & \cdots & y_{M+1} \\
\vdots & \ddots & \vdots & \ddots & \vdots \\
y_{N-1} & \cdots & y_N & \cdots & y_{N+M} \\
\end{bmatrix}, \quad y = \begin{bmatrix}
y_1 \\
y_2 \\
\vdots \\
y_N \\
\end{bmatrix}.
\]

(2.2.8)

Then the AR model coefficients satisfy

\[
\hat{a} = (X^T X + D^T D)^{-1} X^T y,
\]

(2.2.9)

and the residual sum of squares is

\[
S(\lambda^2, \nu^2) = y^T y - \hat{a}^T (X^T X + D^T D) \hat{a}.
\]

(2.2.10)

The likelihood of the hyperparameter model is

\[
L(\lambda^2, \nu^2, \sigma^2) = \left( \frac{1}{2\pi\sigma^2} \right)^{N/2} |D^T D|^{1/2} |X^T X + D^T D|^{-1/2} \exp\left( -\frac{1}{2\sigma^2} S(\lambda^2, \nu^2) \right).
\]

(2.2.11)

Given \( \lambda^2 \) and \( \nu^2 \), the maximum likelihood estimate of \( \sigma^2 \) is, \( \hat{\sigma}^2 = S(\lambda^2, \nu^2) / N \). The ML estimates of \( \lambda^2 \) and \( \nu^2 \) are obtained by minimizing

\[
-2 \log L(y | \lambda^2, \nu^2, \sigma^2) = N \log 2 \pi \sigma^2 + \log |D^T D| - \log |X^T X + D^T D| + N.
\]

(2.2.12)

with respect to \( \lambda^2 \) and \( \nu^2 \). Computation of the likelihood over a discrete \( k, \lambda^2, \nu^2 \) parameter grid
and searching over the resulting discrete likelihood-hyperparameter space for the minimum of $-2 \log$ likelihood yields the desired smoothness priors long AR model.

The frequency domain smoothness priors constraint used here has an interpretation as a constraint on the smoothness of the whitening filter of the AR model. (The 0th derivative has an energy constraint interpretation.) An important facet of our computations is that they are computationally tractable. That allows us to remain within the framework of the general linear model.

An Example, Analysis of Canadian Lynx Data

The data example discussed here is the analysis of the Canadian Lynx data ($n=114$). Other examples are shown in Kitagawa and Gersch (1985a).

The AR model order was set to 20 and up to the fourth order smoothness prior constraint was tried. The hyperparameters $\lambda$ and $\nu$ were searched over the discrete values $\lambda = 2^{j-34} \sigma_0, j=1,\ldots,10$ where $\sigma_0^2$ is the sample variance of the data and $\nu = i^2 \sigma_0, i=0,1,\ldots,4$, for each value of the order of the smoothness prior constraint.

The overall best model was $k=1, \nu = 0, \lambda = 0.173$. The Bayesian estimate of the spectrum is shown in Figure 2a. For comparison, AR models of order up to 20 were fitted by the least squares method. The AIC best order was 11. The AIC criterion-AR modeled spectrum is shown in Figure 2b. In the Bayesian model spectral estimate, the peaks at the high frequencies are significantly reduced compared to the AR model model spectrum estimate, while the ones in the lower frequencies remain unchanged. Figure 2c shows the superimposed estimated spectra obtained from AR models with different model orders. From Figures 2b, 2c we see that the shape of the two right-most peaks of Figure 2.1B vary considerably with model order. Thus they are not estimated reliably by fixed order models. That is typical of the problem of estimating spectral density by fixed order parametric models. If the model fitted to the data is not in the class of models which generated the data, the model fitting is only approximate. The selection of the best non Bayesian parametric model ignores the evidence, in the Bayesian sense, for other parametric models when in
fact it should be taken into account. The suppression of those peaks by the smoothness priors-long
AR model method, shown in Figure 2a, therefore seems quite reasonable.

2.3 DISCUSSION

The variation of the behavior of the solution, from very rough to very smooth, in the trend
estimation problem under the smoothness priors constraints for different values of the hyperparam-
eters, is characteristic of the profound effect of the hyperparameter. The log likelihood of the
hyperparameter versus the hyperparameter changes gradually in the vicinity of the maximum log
likelihood. That fact permits a discrete likelihood-hyperparameter search procedure to be used in
conjunction with a Householder transformation algorithm to realize a reasonable computational
procedure.

The stochastically perturbed difference equation constraint in the trend estimation problem
is a time domain smoothness priors constraint. Akaike (1980), considered other time domain
smoothness constraint problems including the Shiller distributed lag problem and the seasonal
adjustment of time series. Ishiguro et al. (1981), used time domain smoothness constraints and
fixed effects regression in the analysis of earth tide data. The Householder transformation least
squares algorithm FORTRAN programs BAYSEA and BAYTAP-G in TIMSAC-84 (Akaike et al.
1985), are suitable for seasonal adjustment and tidal analyses respectively.

Akaike (1979) illustrated a frequency domain smoothness prior for the solution of the Shiller
(1973) impulse response estimation problem. We used frequency domain smoothness priors here
for spectral density estimation. Gersch and Kitagawa (1984) is an application of the frequency
domain-smoothness priors method to transfer function estimation. Our smoothness priors
method is particularly suited for the situation in which only a short span of data is available for
analysis. In that situation, the results of conventional parametric model analysis methods are par-
ticularly sensitive to the choice of model order. We circumvent that problem, with the Bayesian
smoothness priors method, by tending to fit models that are "too long". The model parameters
are specified as the solution to a constrained least squares problem in which the constraints are expressed in the frequency domain. The likelihood of the hyperparameters is readily computable in a least squares framework with the frequency domain priors.

The goodness of the choice of the frequency domain smoothness priors can be appraised by evaluating its performance in various conceivable situations. The smoothness priors-long AR model gives reasonable results in the analysis of the real Lynx data in comparison with the minimum AIC-AR model method. Kitagawa and Gersch (1985a) show smoothness priors-long AR model results that were superior to the minimum AIC-AR model method in a simulated-two sine waves in noise case, when the data actually corresponds to an ARMA model. This flexibility of performance is what is desired from a Bayesian model.

Also in Kitagawa and Gersch (1985a), a Monte Carlo study of the expected entropy experiment was done to appraise the performance of the smoothness priors-long AR model for spectral estimation against performance of parametric AR models whose order was determined by Akaike's AIC criterion. The smoothness priors-long AR model method was superior to the minimum AIC-AR model method in the two simulation model cases studied. In one case, the simulation model was in the AR model class. In the other case, the simulation model was not in the AR model class. Thus, the example shown here and the Monte Carlo study reported in Kitagwa and Gersch (1985a) are evidence to support the soundness of our empirical frequency domain smoothness priors approach.
3 STATE SPACE GAUSSIAN SMOOTHNESS PRIORS MODELING

A state space modeling approach for the linear model with Gaussian system and observation noise that is the equivalent of the least squares computational approach to smoothness priors modeling, was shown in Brotherton and Gersch (1981) and Kitagawa (1981). The state space smoothness priors modeling method was applied to the modeling of nonstationary mean and nonstationary covariance time series, Gersch and Kitagawa (1983a, 1985) and Kitagawa and Gersch (1984, 1985b).

In the modeling of nonstationary time series discussed below, there tends to be more parameters than data. In that case, attempts to fit the parameters by least squares or any other ordinary means will yield poor parameter estimates. The smoothness priors permit the model parameters to be expressed implicitly as the solution of zero-mean unknown variance stochastically perturbed difference equations. The variances are hyperparameters of the prior distribution of the parameters. One interpretation of the role of the smoothness priors is that they permit a realization of a computational procedure to estimate the model parameters.

In this section, computational procedures for the modeling of nonstationary mean and nonstationary covariance time series are discussed that are variations of the procedures discussed in our previous papers. Examples are shown in Section 3.4. A discussion of other problems treated by the smoothness priors-state space-linear-Gaussian model and comments appear in Section 3.5. Kalman filter, prediction and smoothing formulas and computation of the likelihood of the linear Gaussian model are shown in Section 3.2.

3.1 NONSTATIONARY MEAN SMOOTHNESS PRIORS STATE SPACE MODELING

Time series with trend and seasonal components occur for example in meteorological, oceanographic and econometric studies. Here we consider a complex nonstationary mean time series problem motivated by economic time series considerations. The economic time series nonstation-
ary mean can be decomposed into a trend $t_n$, a globally stationary component $v_n$, a seasonal component $\varepsilon_n$, a trading day factor $d_n$ and an observation noise component $\epsilon_n$,

$$y_n = t_n + \varepsilon_n + v_n + d_n + \epsilon_n. \tag{3.1.1}$$

Each of the aforementioned components can be modeled as a stochastically perturbed difference equation. The generic state space model for this decomposition can be expressed by

$$z_n = Fz_{n-1} + Gw_n \tag{3.1.2}$$

$$y_n = H_n z_n + \epsilon_n$$

where $F, G$ and $H_n$ are $M \times M, M \times L$ and $1 \times M$ matrices respectively. $w_n$ and $\epsilon_n$ are each assumed to be zero mean independent normally distributed random variables. $z_n$ is the state vector at time $n$ and $y_n$ is the observation at time $n$. For any particular model of the time series, the matrices $F, G$ and $H_n$ are known and the observations are generated recursively starting from an initial state that is assumed to be normally distributed with mean $z_0$ and covariance matrix $V_0$.

The state space model that includes the local polynomial trend, stationary AR coefficient, trading day effects and observation error components can be written in the orthogonal decomposition form

$$z_n = \begin{bmatrix} F_1 & 0 & 0 & 0 \\ 0 & F_2 & 0 & 0 \\ 0 & 0 & F_3 & 0 \\ 0 & 0 & 0 & F_4 \end{bmatrix} z_{n-1} + \begin{bmatrix} G_1 & 0 & 0 & 0 \\ 0 & G_2 & 0 & 0 \\ 0 & 0 & G_3 & 0 \\ 0 & 0 & 0 & G_4 \end{bmatrix} w_n \tag{3.1.3}$$

$$y_n = [H_1 \ H_2 \ H_3 \ H_{4,n}] z_n + \epsilon_n.$$

The component models $(F_j, G_j, H_j)$ in order, $(j=1,\ldots,4)$ represent the trend, stationary AR, seasonal and trading day effects components respectively. Some of the particular trend, AR, seasonal and trading day difference equation constraints that we have employed and that have representations in the $(F_j, G_j, H_j)$ matrices in (3.1.3) are shown below.
The trend component $t_n$ satisfies a kth order stochastically perturbed difference equation

$$\nabla^k t_n = w_{1,n}$$

(3.1.4)

where $w_{1,n}$ is an i.i.d. sequence with $w_{1,n} \sim N(0, \tau^2)$. (See (2.1.12).)

The stationary AR component $v_n$ is assumed to satisfy an AR model of order $p$. That is given by

$$v_n = a_1 v_{n-1} + \cdots + a_p v_{n-p} + w_{2,n}.$$  

(3.1.5)

In (3.1.5) $w_{2,n}$ is an i.i.d. sequence with $w_{2,n} \sim N(0, \tau^2)$. The seasonal component of the period L difference equation is

$$s_n = -s_{n-1} - s_{n-2} - \cdots - s_{n-L+1} + w_{3,n}.$$  

(3.1.6)

In (3.1.6), $w_{3,n}$ is an i.i.d sequence with $w_{3,n} \sim N(0, \tau^2)$.

The trading day effect model is

$$d_n = \beta_{1,n} d_{1,n} + \cdots + \beta_{6,n} d_{6,n}$$

(3.1.7)

where $\beta_{i,n}$ denotes the trading-day effect factor and $d_{i,n}$ corresponds to the number of the ith day of the week at time $n$. Implicit in (3.1.7) is the constraint \( \sum_{i=1}^{7} \beta_{i,n} = 0 \). There is no stochastic component in (3.1.7).

For a general model including local polynomial trend, AR component trend, local seasonal component and trading day effect components, the state or system noise vector and observation noise $\epsilon_n$ are assumed to i.i.d. with zero mean and diagonal covariance matrix

$$
\begin{pmatrix}
  w_n \\
  \epsilon_n
\end{pmatrix} \sim N
\begin{pmatrix}
  0 \\
  0 \\
  0 \\
  0 \\
  0 \\
  0 \\
  0
\end{pmatrix},
\begin{pmatrix}
  \tau^2 & 0 & 0 & 0 \\
  0 & \tau^2 & 0 & 0 \\
  0 & 0 & \tau^2 & 0 \\
  0 & 0 & 0 & \sigma^2
\end{pmatrix}
$$

(3.1.8)

An example of a state space model that incorporates each of the components with trend order 2,
AR model order 2 and seasonal component with period L is,

\[
\begin{bmatrix}
  t_n \\
  t_{n-1} \\
  \vdots \\
  t_0 \\
  s_n \\
  s_{n-1} \\
  \vdots \\
  s_0 \\
  \beta_{t,n} \\
  \beta_{s,n}
\end{bmatrix}
= 
\begin{bmatrix}
  2 & -1 & 0 & 0 & 0 & \cdots & 0 & 0 & 0 \\
  1 & 0 & 0 & 0 & 0 & \cdots & 0 & 0 & 0 \\
  \vdots \\
  0 & 0 & a_1 & a_2 & 0 & \cdots & 0 & 0 & 0 \\
  \vdots \\
  0 & 0 & 1 & 0 & 0 & \cdots & 0 & 0 & 0 \\
  0 & 0 & 0 & 0 & -1 & \cdots & -1 & 0 & 0 \\
  \vdots \\
  0 & 0 & 0 & 0 & 0 & \cdots & 1 & 0 & 0 \\
  0 & 0 & 0 & 0 & 0 & \cdots & 0 & 0 & 0 \\
  0 & 0 & 0 & 0 & 0 & \cdots & 0 & 0 & 0 \\
  \beta_{t,n} & 0 & 0 & 0 & 0 & \cdots & 0 & 0 & 0 \\
  \beta_{s,n} & 0 & 0 & 0 & 0 & \cdots & 0 & 0 & 0 \\
\end{bmatrix}
\begin{bmatrix}
  z_{n-1} \\
  u_n
\end{bmatrix}
\]

\[y_n = \begin{bmatrix} 1 & 0 & 1 & 0 & 1 & \cdots & 0 & d_{t,n} & \cdots & d_{s,n} \end{bmatrix} z_n + \epsilon_n.\]  

The smoothness priors problem that includes all of the components in the decomposition identified above corresponds to the maximization of

\[
\exp\left\{-\frac{1}{2\sigma^2} \sum_{n=1}^{N} |y_n - t_n - s_n - d_n|^2 \right\} 
\cdot 
\exp\left\{-\frac{\tau^2}{2\sigma^2} \sum_{n=1}^{N} |\epsilon^2 t_n|^2 \right\} 
\cdot 
\exp\left\{-\frac{\sigma^2}{2\sigma^2} \sum_{i=1}^{2} |a_i| \epsilon_{s,n-i} \right\} 
\cdot 
\exp\left\{-\frac{\tau^2}{2\sigma^2} \sum_{n=1}^{N} \sum_{i=0}^{L-1} |s_n - i|^2 \right\}. 
\]

The first term in (3.1.10) corresponds to the conditional data distribution. The remaining terms in (3.1.10), in order, correspond to the priors on the trend, the globally stochastic component and the seasonal component.

The role of the hyperparameters \(\tau^2\) and \(\sigma^2\) as a measure of the uncertainty in the belief of the priors is clear from (3.1.10). Relatively small \(\tau^2\) (\(\sigma^2\)) imply relatively wiggly trend (seasonal) components. Relatively large \(\tau^2\) (\(\sigma^2\)) imply relatively smooth trend (seasonal) components. Correspondingly, the ratio of \(\tau^2/\sigma^2\), \(j=1\) or 3, can be interpreted as signal-to-noise ratios. (The
value of $\sigma^2$ in (3.1.10) is essentially estimated free of computational cost in the Kalman filter algorithm.)

3.2 RECURSIVE ESTIMATION OF STATE AND LIKELIHOOD COMPUTATION

Let a state space model be given by

\[
\begin{align*}
    x_n &= F_n x_{n-1} + G_n w_n \\
    y_n &= H_n x_n + \epsilon_n,
\end{align*}
\]

(3.2.1)

where $w_n \sim N(0, Q_n)$ and $\epsilon_n \sim N(0, R_n)$. Given the observations $y_1, \ldots, y_N$ and the initial conditions $x_{0|0}, V_{0|0}$, the one-step-ahead predictor and the filter are obtained from the Kalman filter algorithm:

**Time Update (Prediction)**

\[
\begin{align*}
    x_{n|n-1} &= F_n x_{n-1|n-1} \\
    V_{n|n-1} &= F_n V_{n-1|n-1} F_n^T + G_n Q_n G_n^T.
\end{align*}
\]

(3.2.2)

**Observation Update (Filtering)**

\[
\begin{align*}
    K_n &= V_{n|n-1} H_n^T H_n V_{n|n-1} H_n^T + R_n \]^{-1} \\
    x_{n|n} &= x_{n|n-1} + K_n [y_n - H_n x_{n|n-1}] \\
    V_{n|n} &= [I - K_n H_n] V_{n|n-1}.
\end{align*}
\]

(3.2.3)

Using these estimates, the smoothed value of the state $x_n$ given the entire observation set, $y_1, \ldots, y_N$, is obtained by the fixed interval smoothing algorithm, (Anderson and Moore 1979),

\[
\begin{align*}
    A_n &= V_{n|n} F_n^T V_{n+1|n}^{-1} \\
    x_{n|N} &= x_{n|n} + A_n [x_{n+1|N} - x_{n+1|n}] \\
    V_{n|N} &= V_{n|n} + A_n [V_{n+1|N} - V_{n+1|n}] A_n^T.
\end{align*}
\]

(3.2.4)

The state space representation and the Kalman filter yield an efficient algorithm for the likelihood of a time series model. The joint distribution of $y_1, \ldots, y_N$ is,
\[ f(y_1, \ldots, y_N) = \prod_{n=1}^{N} f(y_n | y_1, \ldots, y_{n-1}), \] (3.2.5)

with

\[ f(y_n | y_1, \ldots, y_{n-1}) = (2\pi v_n)^{-1/2} \exp \left\{ -\frac{1}{2v_n} (y_n - H_n x_{n|n-1})^2 \right\}, \] (3.2.6)

\[ v_n = H_n V_{n|n-1} H_n^T + R_n. \]

Then, the log likelihood, \( l \), of the model is obtained by

\[ l = -\frac{1}{2} \left[ N \log 2\pi + \sum_{n=1}^{N} \log v_n + \sum_{n=1}^{N} \frac{-1}{2v_n} (y_n - H_n x_{n|n-1})^2 \right], \] (3.2.7)

The maximum likelihood estimate of the model parameters are obtained by maximizing (3.2.7) with respect to those parameters. The AIC is defined by

\[ AIC = -2(\text{maximum log-likelihood}) + 2(\text{number of parameters}) \] (3.2.8)

Alternative models for time series might be models with and without trading day effects or models with and without AR component effects. In each case, when we consider alternative models, the model with the smallest value of the AIC statistic is selected as the AIC best model.

In fitting a stationary model, we can utilize the theoretical mean and the theoretical covariance of the state vector as the initial values \( x_{0|0} \) and \( V_{0|0} \). In the nonstationary case we consider the initial vector, \( x_{0|0} \), as an unknown parameter and estimate it by using the entire set of data.

The log likelihood obtained by estimating the initial state vector is a natural estimate of the expected log likelihood of the predictive distribution (Akaike 1980b, and Gersh and Kitagawa 1983a).

3.3 NONSTATIONARY COVARIANCE MODELING

Here, time series with nonstationary covariances are modeled by a time varying autoregressive (AR) model with smoothness constraints on the AR parameters. Time varying AR coefficient models have been a topic of research for some time. For example, see Whittle (1965), Kozin
(1977) and Nicholls and Quinn (1985) and the extensive references therein, particularly to the use of random coefficients in econometric modeling. Earlier, in engineering applications, the modeling of nonstationary covariance time series was done by fitting locally stationary models, and by orthogonal polynomial expansions of AR coefficient models. Astrom and Wittenmark (1973, Theorem 5) express a time varying AR coefficient model that includes the possibility of random AR coefficients. Bohlin (1976) is an early application of the analysis of time series models with time-dependent coefficients. The concept of the likelihood of the Bayesian model or of hyperparameters or anything related to a smoothness prior do not appear in the earlier papers. Those are key concepts here. Kitagawa (1983) is a precedent to the material discussed in this section.

The problem in modeling nonstationary covariance time series is to achieve an efficient parameterization to capture the local and global statistical relationships in the time series. That objective is achieved here by imposing smoothness priors constraints in the form of stochastically perturbed difference equations on the evolution of the AR coefficients. The variances of the white noise stochastic perturbations are the hyperparameters of the the AR coefficient distribution. The difference equations are imbedded into a state-space representation. A relatively large AR model order is chosen, the AR coefficients at each time instant are also smoothed using the frequency domain differential constraints on AR coefficients, as in the smoothness priors-long AR model for spectral estimation, Section 2.2. For each order of the differential constraint, the Kalman filter yields the likelihood of the hyperparameters. The smoothed estimate of the nonstationary innovations series variance is also computed. That is used in the computation of an instantaneous spectral density which is defined in terms of the instantaneous AR model coefficients and the innovations series variance.

The Time Varying AR Coefficient Model

A time varying AR coefficient model is given by
\[ y_n = \sum_{i=1}^{m} a_{i,n} y_{n-i} + \epsilon_n. \quad (3.3.1) \]

In (3.3.1), the coefficients \( a_{i,n} \) are assumed to change "gradually" with time and \( \epsilon_n \) is assumed to be a normally distributed white noise sequence with variance \( \sigma^2 \). Since there are \( m \times N \) AR coefficients in the model in (3.3.1), an attempt to fit the parameters by least squares or any other ordinary means to the \( N \) observations \( y_1, \ldots, y_N \), will yield poor parameter estimates. We consider the unknown AR coefficients to be random variables and impose stochastic constraints on those coefficients. Those constraints define a Gaussian smoothness prior distribution on the time history of the AR coefficients and on the spectrum at each time instant. A simple and useful model for a time varying AR coefficient model is obtained by the stochastically perturbed difference equation constraint model

\[ \nabla^k a_{i,n} = \delta_{i,n}, \quad i=1,\ldots,m. \quad (3.3.2) \]

For convenience, in (3.3.2) \( \delta_{i,n} \) is assumed to be a zero-mean Gaussian white noise sequence with variance \( \tau_i^2 \) independent of \( i \) and \( n \). That is, \( \tau_i^2 = \tau_i^2, \quad i=1,\ldots,m. \)

The smoothness priors constraints on the AR coefficients mitigate the problem of overparameterization by permitting the AR coefficients to be expressed as the solution of the constrained least squares problem

\[ \sum_{n=1}^{N} \left[ y_n - \sum_{i=1}^{m} a_{i,n} y_{n-i} \right]^2 + \tau^2 \sum_{n=1}^{N} \sum_{i=1}^{m} |\nabla^k a_{i,n}|^2 + \lambda^2 \sum_{n=1}^{N} \sum_{i=1}^{m} a_{i,n}^2 + \nu^2 \sum_{n=1}^{N} \sum_{i=1}^{m} a_{i,n}^2. \quad (3.3.3) \]

In (3.3.3) \( m \) and \( k_1,k_2 \) are assumed known and \( \tau^2,\lambda^2,\nu^2 \) are the tradeoff parameters which balance the infidelity of the model to the data and the infidelity of the model to the smoothness constraints.

Similar to the analysis in Section (2.1), (3.3.3) yields a Bayesian interpretation of the least squares problem. Multiply (3.3.3) by \(-1/2\sigma^2\) and exponentiate. Then, to within a constant term,

\[ \exp \left( -\frac{\tau^2}{2\sigma^2} \sum_{n=1}^{N} \sum_{i=1}^{m} |\nabla^k a_{i,n}|^2 \right) \exp \left( -\frac{\lambda^2}{2\sigma^2} \sum_{i=1}^{m} a_{i,n}^2 \right) \exp \left( -\frac{\nu^2}{2\sigma^2} \sum_{i=1}^{m} a_{i,n}^2 \right) \quad (3.3.4) \]
expresses the product of the prior distribution on the smoothness of the spectrum and the prior distribution on the smoothness of the AR parameters. The tradeoff parameters $r^2, \lambda^2, \nu^2$ are the hyperparameters of the prior distribution. As in the development in Section (2.1), the product of the conditional data distribution, (proportional to the leftmost term in (3.3.3)), and the prior distribution in (3.3.4) yields the posterior distribution for the AR parameters. As in (2.1.8), integration of the posterior distribution for the AR parameters yields the likelihood for the smoothness tradeoff parameters.

State Space Time Varying AR Coefficient Model

Define the $km\times1$ state vector at time $n$ to be $x_n = (a_{1,n}, \ldots, a_{m,n}, \ldots, a_{1,n-k+1}, \ldots, a_{m,n-k+1})^T$. The state space time varying AR coefficient model is

$$x_n = Fx_{n-1} + Gw_n$$

$$y_n = H_n x_n + v_n.$$  \hfill (3.3.5)

In (3.3.5), $H_n$ is a $(m+1)\times km$ observation matrix, $w_n$ is the $m$ vector, $w_n = (\xi_{1,n}, \ldots, \xi_{m,n})^T$ and the $m+1$ vector $v_n = (\epsilon_n, \xi_n)^T$, is defined in (3.3.7). For the difference equation orders $k_1 = 1$ and $k_2 = 2$, the matrices $F, G, \text{and } H$ are

$$k=1: F = (I_m), \quad G = (I_m), \quad H_n = H_1,n = [y_{n-1}, \ldots, y_{n-m}]$$

$$k=2: F = \begin{bmatrix} 2I_m & -I_m \\ I_m & 0 \end{bmatrix}, \quad G = \begin{bmatrix} I_m \\ 0 \end{bmatrix}, \quad H_n = H_{2,n} = [H_{1,n} \ 0 \ldots \ 0].$$ \hfill (3.3.6)

The input process noise $w_n$, the observation noise $\epsilon_n$ and the spectrum smoothness priors noise form the $2km+1$ vector $(w_n^T, \epsilon_n, \xi_n)^T$ that is assumed to be normally distributed and independent with mean and covariance matrix,

$$\begin{pmatrix} w_n \\ \epsilon_n \\ \xi_n \end{pmatrix} \sim N \left( \begin{bmatrix} 0 \\ 0 \end{bmatrix}, \begin{bmatrix} Q & 0 & 0 \\ 0 & 0 & 0 \\ 0 & 0 & S \end{bmatrix} \right).$$ \hfill (3.3.7)

In (3.3.7), the $m \times m$ diagonal matrix $Q$ has the the element $1/r^2$ on the diagonal and for $k=2$, (in
(3.3.3)), the $m \times m$ diagonal matrix $S$ has diagonal elements $1/(\lambda^2 + \nu^2, \ldots, m^4 \lambda^2 + \nu^2)$, (see (2.2.6)).

For a fixed difference order $k_1$, the best fit of the state space smoothness priors constraints-time varying AR coefficient model to the data $y_1, \ldots, y_N$, is the one for which the likelihood of the hyperparameters $r^2, \lambda^2, \nu^2$, are maximized. The likelihood is computed using the recursive formulas indicated in Section 3.2.

The Instantaneous Variance And The Instantaneous Spectrum

In many practical data analysis situations, the relatively fast wiggles of a nonstationary covariance time series appears to be modulated by a relatively slowly changing envelope function. The envelope function has an interpretation as a change of scale of the time varying AR coefficient model or equivalently as the smoothed (trend) value of the instantaneous variance, (Section 3.1). A key idea in that modeling is to find a variance stabilizing transformation of the innovations that yields the instantaneous trend in an additive zero-mean constant variance observation noise, Wahba (1980). Let $s_n, n=1,\ldots, N$ be a realization of a zero-mean normally distributed white noise with unknown variance $\sigma^2_n$. Then, if $\sigma^2_n = \sigma^2_{n-1}, \chi^2_m = |\xi^2_{m-1} + \xi^2_m|/2$, is an independent sequence of chi-square random variables with two degrees of freedom. From Wahba (1980), the transformation $t_m = \log(\chi^2_m) + \gamma$, where $\gamma = 0.57721$ is the Euler constant, leaves the independent random variable $t_m$ with a distribution that is almost normal with the moments $E[t_m] = \log \sigma^2_m, \ Var[t_m] = \pi^2/6$. That idea is exploited here.

Consider a second order difference equation constraint on the log variance defined by

$$\nabla^2 t_m = w_m,$$  \hspace{1cm} (3.3.8)

In (3.3.8) $\{w_m\}$ is an independent zero-mean normally distributed sequence with unknown variance $r^2$. Define a state vector by $z_m = [t_m, t_{m-1}]^T$. Then in state space form the constraint model in (3.3.8) and the observation model are given by
\[
\begin{bmatrix}
-2 & 1 \\
1 & 0
\end{bmatrix}
\begin{bmatrix}
\xi_m \\
\xi_{m-1}
\end{bmatrix}
+
\begin{bmatrix}
1 \\
0
\end{bmatrix}
\omega_m
\]

\[
\eta_m = \begin{bmatrix} 1 & 0 \end{bmatrix} \xi_m + \epsilon_m.
\]

Application of the Kalman filter, prediction and smoothing algorithms described earlier yield the smooth value \(t_{m|N}\), the logarithm of the smoothed estimate of the changing variance. Our estimate of the changing variance is \(\sigma^2_{2m|N} = \sigma^2_{2m-1|N} = \exp(t_{m|N} + \gamma)\).

Motivated by earlier work on spectrum estimation, we define the instantaneous spectrum of a time varying coefficient AR process by

\[
S_n(f) = \frac{\sigma_n^2}{1 - \sum_{k=1}^{m} a_k \exp(-2\pi i kf)}; \quad -1/2 \leq f \leq 1/2.
\]

The value of the instantaneous spectrum is obtained by substituting the smoothed estimates of the time varying AR coefficients and the smoothed estimate of the innovations variance \(\sigma_n^2\) into (3.3.10).

### 3.4 EXAMPLES

#### A Nonstationary Mean Time Series Example

The RSWomen series of the Bureau of the Census data, \(\text{Zellner 1983}\), is analyzed here. The series consists of the retail sales of women's apparel, reported in millions of dollars. The sales for each month are affected by the number of times each day of the week occurs, the trading-day effect, because buying behavior differs for each day of the week. (The sales are also affected by holidays.) We are interested in determining whether or not it is appropriate to include a trading day effect in the model and whether or not the globally stochastic AR component should be included in the model. The AIC statistic is used to determine the best of the alternative models.

The computations were done using the DECOMP.FORT program in TIMSAC-84 (Akaike et al. 1985). The model was fitted to the data \(y_1, \ldots, y_{129}\), 24 data points were withheld. The AIC's for
AR model orders $p=0,1,2,3$, (as in (3.1.5)), respectively for the non trading day effect and trading day effect models are $(111.51,96.96,98.80,98.65)$ and $(88.07,68.34,67.23,68.91)$. An interpretation of those results is that models with an AR component, $p \neq 0$, are superior to models without AR components both with and without trading day effects and that the AIC best model, $(AIC=67.23)$, is the trading day effect model with AR model order $p=2$. Figures 3a-e show selected computational responses for the non-AR component-trading day effect model. Figures 3f-j show selected computational results for the AR component-trading day effects model. The seasonal components, residual noise and trading day effects and seasonal plus trading effects are quite similar in appearance for both models.

Several aspects of the modeling results are noteworthy. The trend of the trend plus AR component model is smoother than the trend-non AR component model and the trend plus AR component is almost indistinguishable from the trend in the non AR component model. Also, the seasonal component is very regular whereas the seasonal plus trading day component reveals the expected slight irregularities.

A important property of the AIC best trend plus seasonal plus AR component model, instead of the trend plus seasonal model, can be seen in the (out-of-sample) forecasts for these models as shown in Figures 3e and 3j. In those illustrations we show the true series, the forecasted series and plus and minus one sigma of the forecast random variable. The plus or minus one sigma prediction intervals of the trend plus AR plus seasonal plus trading day components model is much tighter than the same quantity for the non AR component model. The increase in prediction variance per step in increasing horizon forecast, grows in accordance with the variance component terms in the matrix, in (3.1.8). The variance of the (wiggly) trend term in the non AR component model, is larger than the sum of variance terms of the (smooth) trend and AR component in the AR component model. That larger variance is reflected into the larger one sigma prediction interval of the non AR component model.
These results illustrate the flexibility of the decomposition of the nonstationary mean concept via smoothness priors modeling and the importance of the role of the AIC in selecting the best of alternative models.

Time Varying AR Coefficient Modeling, Nonstationary Covariance Time Series

The computations were realized using the TVCAR.FORT program in TIMSAC-84 (Akaike et al 1985). Figure 4a shows a seismic data event, $y_1, ..., y_N$. The stochastic "background noise", P wave (the first abrupt change in the signal) and the S wave (the second abrupt change in the signal), are clearly discernable. Figures 4bd,f are graphs of computational results from the time varying AR coefficient model described in Section 3.3. Respectively they show the $\log ((y_m^2 + y_{m-1}^2), m=1, ..., N/2$ "unsmoothed envelope" data and the superimposed estimate of the envelope (changing variance), the evolution of the instantaneous power spectral density and the evolution of the partial correlation coefficients (parcors) of the fitted, AR order $m = 8$ time varying AR coefficient model. Figures 4c,e,g show the corresponding computational results from an "intervention analysis" model that is part of TVCAR.FORT.

The smoothed envelope for the non-intervention model is in fact quite smooth. Similarly the instantaneous spectrum and partial correlation coefficients (parcors) reflect the smooth transition in the time varying AR coefficients model, (3.3.2). Two visual inspection determined "outlier" events occur at $n=635$ and $n=1030$. They correspond to the arrival times of the P and S waves and are identified to the program, by human intervention, as large observation variance events. The large observation variance relaxes the priors constraint and permits the AR coefficients and subsequently derived quantities to change abruptly at those instants. The "validity" of this intervention type analysis is suggested by comparison of the results of the intervention type and non-intervention type analysis. The properties of the latter "drift" toward the former. The abrupt changes in the appearance of the envelope function and in the instantaneous spectra and parcors correspond to the physical interpretation that the P waves and S waves are different sources of
energy at the observing seismometer.

3.5 COMMENTS, DISCUSSION

The paper by Akaike (1980) motivated our work in smoothness priors. In Akaike (1980), computations were done by a Householder transformation least squares algorithm of computational complexity $O(N^2)$. Brotherton and Gersch (1981) and Kitagwa (1981) demonstrated an equivalent state space modeling approach for the linear model with Gaussian system and observation noise. In the vicinity of the maximized likelihood, the likelihood is a rather flat function of the hyperparameters. This fact permits a relatively coarse grid discrete hyperparameter-likelihood search procedure to determine the values of the hyperparameters that tend to maximize the likelihood. Such a procedure preserves the $O(N)$ computational complexity inherent in the Kalman filter computations. A computational complexity of $O(N)$ version of that method was subsequently applied to a variety of nonstationary mean and nonstationary covariance time series modeling problems, (Gersch and Kitagwa 1983,1985, Kitagawa and Gersch 1984, 1985b). Variations of the procedures in those papers expressed in computer programs DECOMP.FORT and TVCAR.FORT (TIMSAC-84, Akaike et al. 1985), yielded the computational results shown here.

Potentially many more combinations and extensions of the models shown here for the modeling of nonstationary mean and nonstationary covariance time series by smoothness priors methods, are possible. For example, a generalization of the regression on trading days components, in the nonstationary mean-decomposition of time series modeling, could take into account constant coefficient and/or time varying coefficient regressssion on other time series. A time varying partial AR coefficients variation of the present time varying AR coefficient model for the modeling of nonstationary covariance time series, has already been implemented. Another potential variation on the time varying AR coefficient model would be to estimate the full nondiagonal $m \times m$ system noise covariance matrix (the matrix of hyperparameters). Gersch and Kitagawa (1985), an application of the time varying AR coefficient model, includes computation of the time
varying covariance function. That computation is useful to permit computation of the mean square response to nonstationary excitation of building structures to single realizations of seismic event data.

Some other linear model-Gaussian disturbances- state space smoothness priors models have been implemented. Kitagawa and Takanami (1985) show a smoothness priors modeling method for the extraction of seismic signals from correlated background noise. The smoothness priors innovation in that work is the implementation of a non constant or time varying hyperparameter. That hyperparameter achieves a time varying balance of the tradeoff between the variances of the seismic signal and the background noise. The modeling of continuous model time series with discrete time observations is another domain where smoothness priors state space modeling has been exhibited. Kitagawa (1984) includes a smoothness priors variation of the Jones (1980) continuous time AR process-discrete time observations modeling. The application in Kitagawa (1984) is to irregularly spaced or missing data time series modeling.
4. STATE SPACE NON GAUSSIAN MODELING OF NONSTATIONARY MEAN TIME SERIES

A non-Gaussian state space approach to the modeling of nonstationary time series is shown. Neither the system noise nor the observation noise need be Gaussian. Recursive formulas for the prediction, filtering and smoothing of the state are given. A numerical method, based on a piecewise linear approximation to the density functions for realizing these formulas, is also given. The merits and potential wide applicability of this approach to non-Gaussian modeling are illustrated by some numerical examples. Extension of this method to the state space modeling of nonlinear systems is straightforward.

Earlier in this chapter we demonstrated the wide range of applicability of the linear model with Gaussian system and observation noise. There are numerous problems for which Gaussian modeling is inadequate. For example, the problem of trend estimation becomes difficult when the trend has discontinuities as well as smooth changes and when there are observation outliers. A simple linear Gaussian model with small process noise variance does not track jumps or discontinuities very well. A model with large process noise variance will respond to sudden changes in the trend but it will also be inappropriately wiggley where the trend is quite smooth. The treatment of such trend discontinuities with the included possibility of observation outliers in the linear Gaussian model framework requires a complicated model. Heavy tailed distributions for process and observation noise can cope with these problems with a simple model. Also, smoothing problems in which there is a time varying variance and/or a nonhomogenous binomial or Poisson mean require a non Gaussian system noise model formulation. Similarly nonlinear models such as storage models for river flow and a ship's nonlinear maneuverability require non Gaussian distribution models.

Thus, the development of methods for treating systems with non Gaussian distributions is well motivated. In earlier attempts, systems with non Gaussian distributions were approximated by the use of extended Kalman filters, sums of Gaussian distributions, by Edgeworth or Gram-
Charlier expansions etc... (See Alspach and Sorenson 1972, for example.) Here, we approximate the probability density functions directly by a piecewise linear function. The recursive prediction, filtering and smoothing computation required by the state space modeling are realized by numerical integration. A similar approach was considered by Bucy and Senne (1971) and de Figueiredo and Jan (1971) in the context of nonlinear filtering problems. Such an approach is more feasible now than it was several years ago because of the development and proliferation of fast computational facilities. In Section 4.1 the state space prediction, filtering and smoothing formula aspects of the numerical computations are derived and the computation of the likelihood for the not necessarily Gaussian distribution model are shown. Numerical examples are shown in Section 4.2 and a discussion and comments are in Section 4.3.

4.1 THE NON GAUSSIAN STATE SPACE MODEL

Consider the stationary state space system described by

\[ x_n = F x_{n-1} + G w_{n-1} \] (4.1.1)

\[ y_n = H x_n + \epsilon_n \]

where as before F, G, and H are linear transformations. The independent and independent of each other, but not necessarily Gaussian process and observation noises are \( w_n \) and \( \epsilon_n \) respectively. The initial state vector \( x_0 \) is distributed in accordance with \( p(x_0) \) and the conditional density of the state at time \( n \), given the observations \( (y_1, \ldots, y_m) = Y_m \) is denoted by \( p(x_n | Y_m) \). Then, the recursive formulas for the one step ahead prediction, filtering and smoothing densities are derived as follows:

One step ahead prediction (time update)

\[ p(x_n | Y_{n-1}) = \int_{-\infty}^{\infty} p(x_n, x_{n-1} | Y_{n-1}) dx_{n-1} \] (4.1.2)

\[ = \int_{-\infty}^{\infty} p(x_n | x_{n-1}) p(x_{n-1} | Y_{n-1}) dx_{n-1} \]
Filtering (observation update)
\[
p(z_n \mid Y_n) = p(z_n \mid y_n, Y_{n-1}) = \frac{p(z_n, y_n \mid Y_{n-1})}{p(y_n \mid Y_{n-1})} \tag{4.1.3}
\]
\[
= p(y_n \mid z_n) p(z_n \mid Y_{n-1}) / p(y_n \mid Y_{n-1})
\]
where \(p(z_n \mid z_{n-1})\) is the density of \(z_n\) given the previous state vector \(z_{n-1}\), \(p(y_n \mid z_n)\) is the density of \(y_n\) given \(z_n\) and \(p(y_n \mid Y_{n-1})\) is obtained by \(\int p(y_n \mid z_n) p(z_n \mid Y_{n-1}) \, dz_n\).

Similarly, consider the expression for the joint density of \(z_n\) and \(z_{n+1}\), given the entire observation sequence \(Y_N\),
\[
p(z_n, z_{n+1} \mid Y_N) = p(z_{n+1} \mid Y_N) p(z_n \mid z_{n+1}, Y_n) \tag{4.1.4}
\]
\[
= p(z_{n+1} \mid Y_N) p(z_n, z_{n+1} \mid Y_n) / p(z_{n+1} \mid Y_n)
\]
\[
= p(z_{n+1} \mid Y_N) p(z_n \mid z_{n+1}) p(z_n \mid Y_n) / p(z_{n+1} \mid Y_n)
\]
From (4.1.4) we obtain the formula for smoothing:
\[
p(z_n \mid Y_N) = \int_{-\infty}^{\infty} p(z_n, z_{n+1} \mid Y_N) \, dz_{n+1} \tag{4.1.5}
\]
\[
= p(z_n \mid Y_n) \int_{-\infty}^{\infty} p(z_{n+1} \mid Y_N) p(z_n \mid z_{n+1}) p(z_n \mid Y_n) / p(z_{n+1} \mid Y_n) \, dz_{n+1}.
\]

In the linear Gaussian case, the conditional densities \(p(z_n \mid Y_{n-1})\), \(p(z_n \mid Y_N)\) and \(p(z_n \mid Y_N)\) are characterized by mean vectors and covariance matrices. Correspondingly, (4.1.2), (4.1.3) and (4.1.5) lead to the Kalman filter and the fixed interval smoothing algorithm. In the non Gaussian or nonlinear case however, it is necessary to evaluate the non Gaussian densities explicitly at each step. The algorithms above, (4.1.3)-(4.1.5) can be realized numerically by piecewise linear approximations to the density functions, transformation of densities, convolution of densities and Bayes theorem (product of two densities and normalization). Details of the numerical computations are in Kitagawa (1987).

In general, the non Gaussian model has some unknown parameters. The best choice of the parameters can be found by maximizing the log likelihood defined by
\[ l(\theta) = \log p(y_1, \ldots, y_N) = \sum_{n=1}^{N} \log p(y_n | y_1, \ldots, y_{n-1}) = \sum_{n=1}^{N} \log p(y_n | Y_{n-1}). \] (4.1.6)

The term \( p(y_n | Y_{n-1}) \) appears in (4.1.3) and can be evaluated numerically. If we have several candidate models, including models with different types of system noise or observation noise density functions, we choose the model for which the AIC is minimum.

### 4.2 NUMERICAL EXAMPLES

**Estimation of a Shifting Mean Value**

Consider the data simulated from the following model,

\[ Y_n \sim N(\mu_n, 1) \]

\[ \mu_n = \begin{cases} 
0 & n = 1, \ldots, 100 \\
1 & n = 101, \ldots, 200 \\
-1 & n = 201, \ldots, 300 \\
2 & n = 301, \ldots, 400 
\end{cases} \] (4.2.1)

The data is shown in Figure 5a. The problem is to estimate the abruptly changing mean value function \( \mu_n \).

For this type of data we used the model

\[ \nabla^k t_n = w_n \] (4.2.2)

\[ y_n = t_n + \epsilon_n. \]

As before, \( \nabla \) is the difference operator defined by \( \nabla t_n = t_n - t_{n-1} \) and \( w_n \) and \( \epsilon_n \) are white noise sequences that are not necessarily normally distributed. For simplicity we assume that the difference order \( k \) is one. Equation (4.2.2) is a special form of the state space model, Section (3.1.), with \( z_n = t_n, F = G = H = 1 \). We considered the following model classes:

- **Model(a):** \( w_n \sim a N(0, \tau^2) + (1 - a) N(0, \tau_1^2), \quad \epsilon_n \sim N(0, 1) \) (4.2.3)
- **Model(b):** \( w_n \sim Q(b, \tau^2), \quad \epsilon_n \sim N(0, 1) \)

Model(a) denotes a mixture of Gaussian system noises. In (4.2.3), for Model(a) and Model(b),
\( N(0,r) \) denotes the Gaussian distribution with mean 0 and variance \( r \). In Model(b), \( Q(b,r^2) \)
denotes the distribution of the Pearson system with density \( q(x;b,r^2)\) = \( C(r^2 + z^2)^{-b} \) with \( \frac{1}{2} \leq b \leq \infty \) and \( C = r^{2b-1} \Gamma(b)/\Gamma\left(\frac{1}{2}\right) \). This family links the Cauchy distribution (\( b = 1 \)) and the Gaussian distribution (\( b = \infty \)). In Model(a), \( r^2 \) was arbitrarily set to 4.0, approximately the sample variance of the simulated data. The maximum likelihood estimate of \( r^2 \) for the Gaussian model, Model(a), with \( a = 1.0 \) or equivalently Model(b) with \( b = \infty \), was \( \hat{r}^2 = 0.0429 \). The AIC of the model was 1240.33. For the mixture of Gaussian system noises model, \( \hat{a} = 0.989 \), \( \hat{r}^2 = 0.0000014 \) and \( AIC = 1212.48 \). We tried four Pearson family models: \( b = 0.6, 0.8, 1.0 \) and \( \infty \). \( b = 0.80 \) is the AIC best Pearson family model with \( \hat{r}^2 = 0.000002 \) and \( AIC = 1215.20 \). The AIC best model is the mixture of Gaussian system noises model.

Figure 5b-5d shows the marginal posterior density \( p(z_n | Y_N) \) versus time \( n \) for the Gaussian model, the best Pearson system model and the mixture of Gaussians model respectively. For the Gaussian model, Figure 5b, the densities obtained have identical shape except for the ends of the time interval where the densities become slightly broader. In Figure 5c, the shape of the posterior density varies with time. When the mean value shifts, the density becomes heavy tailed on one side. The Gaussian mixtures model also exhibits the latter behavior.

Figures 5e-5g shows the mean (bold) and \( \pm 1,2,3 \) sigma intervals of the \( p(z_n | Y_N) \) versus \( n \) for the Gaussian model and the median (bold) and corresponding 0.13, 2.27, 15.87, 84.13, 97.73 and 99.87 percentage points of the Pearson system model and the Gaussian mixtures model respectively. For the Gaussian model, the estimated mean value function becomes a wiggly curve and does not reflect the sudden change of the mean value. The estimated median of the Pearson system and the Gaussian mixtures models do capture the sudden change of the signal mean value. The multimodal or skewed distribution and jumps of the mean value are typical of the phenomenon that are seen in non Gaussian modeling.
Estimation Of Changing Variance

The estimation of changing variance was discussed in Section 3.3.3 in the context of fitting a time varying AR coefficient model to a seismic signal. The same idea is exploited here to estimate the changing variance of the same seismic signal with the state space non Gaussian modeling method. A first order difference model for the trend of the log of the sum of the squares of successive observations was used, $\nabla t_m = t_m - t_{m-1} \sim \mathcal{N}(\beta, \tau^2)$, $y_m = t_m + \epsilon_m$ with $m=1,\ldots,N/2$. Two models were considered, the Gaussian system noise-Gaussian observation noise and the Cauchy system noise with an $r(x) = \exp(x - \exp x)$ observation noise model. The corresponding AIC's were 4778.94 and 4222.84. The latter model was the AIC best model. The original seismic wave $y_{n}, n=1,\ldots,N$ and the $\log((y_{2m}^2 + y_{2m-1}^2)/2)$, $m=1,\ldots,N/2$ signals are shown in Figures 6a and 6b respectively. The Gaussian model smoothed mean and $\pm 1,2,3\sigma$ and non Gaussian model smoothed median and corresponding probability point curves are shown in Figures 6c,6d respectively. Those illustrations indicate that the Cauchy system noise model yields better estimates of the smooth mean and abrupt changes of the mean innovations variance than the simple Gaussian system noise model.

Modeling the real data changing variance seismic signal with the state space non Gaussian system noise method automatically yields the location of abrupt changes in the mean of the signal.

4.3 COMMENTS, DISCUSSION

The results shown here were obtained using a simple one dimensional state vector. In principle, it is straightforward to extend the computational formulas to higher dimensional state systems. The resulting increase in the computational burden required to compute the convolution of density functions becomes quite severe. A variety of numerical techniques have been investigated to cope with this problem. Very likely the use of more powerful computers rather than the increase of effort in numerical analysis methods will be more expeditious in the development of the non Gaussian state space smoothness priors method.
Several other problems lend themselves to the application of the one dimensional non Gaussian modeling shown here. Kitagawa (1987) shows an application to the handling of discrete distributions. The time varying mean of a real, nonstationary (nonhomogeneous) binary process, is estimated. Also, smoothing of the log periodogram using a state space model with \( \omega_n \sim \log \chi^2 \) and \( \epsilon_n \sim \text{Cauchy} \) or \( \epsilon_n \sim \text{Gaussian} \) is an alternative to the Gaussian distributions approach in Wahba (1980).

Our procedure also extends quite naturally to the analysis of nonlinear systems. The one-step-ahead prediction formula (4.1.2) and the filtering formula (4.1.3) are applicable even for nonlinear systems.

Time series with nonstationarities, nonlinearities and outliers that have been difficult to analyze by conventional linear Gaussian models can be quite simply analyzed with the non Gaussian model. The computational burden using the non Gaussian filter and smoother is substantial. The development of faster algorithms and the use of faster computers will alleviate this burden.
5. SUMMARY

The ingredients for the smoothness priors analysis of time series are the model, the prescription of the priors, the criterion for goodness of model fit and the computational method.

Initially, smoothness priors modeling of stationary, nonstationary mean, and nonstationary covariance time series was demonstrated in the context of the linear model with Gaussian distributed system noise and with Gaussian distributed observation noise. Both time domain and frequency domain specifications of the prior distribution of the model parameters were considered. A hyperparameter specifies the degree of belief in the prior distribution. The smoothness priors method of analysis derives its unity from the fact that the likelihood of the Bayesian model (the likelihood of the hyperparameter(s)) is the single criterion by which the goodness of fit of the model is determined. The maximization of the likelihood of a small number of hyperparameters permits the modeling of time series with complex structure and a large number of implicitly inferred parameters. When there are alternative candidate smoothness priors models, we use Akaike's AIC to determine the best of alternative models, (the likelihood of the model has a central role in the AIC). Householder transformation least squares and Kalman filter algorithms, were the means for the realization of the smoothness priors time series modeling.

Finally, we demonstrated a state space representation not necessarily Gaussian not necessarily linear model method of smoothness priors modeling. In that method, piecewise-linear approximation to densities and numerical integration computations were employed. Conceptually all of the possible combinations of models and smoothness priors computations could be realized with this method.

The extensive applicability of smoothness priors modeling methods in time series modeling was demonstrated. A large number of other problems that have not been well solved by more traditional time series methods remain to be solved by that method.
REFERENCES


Gersch, W. and Kitagawa, G. (1983), The Prediction of Time Series with Trends and

39


LEGENDS

FIGURE 1. Trend Estimation
a: Truncated Gaussian signal and signal plus noise, b: Signal plus noise plus smoothed trend with a too large hyperparameter, c: Signal plus noise plus smoothed trend with a too small hyperparameter, d: Signal plus noise plus smoothed trend with optimum hyperparameter,

Figure 2. (2.1) Spectral Densities From Canadian Lynx Data Example
a: Spectral density versus frequency, smoothness priors model, b: Spectral density versus frequency, AIC-AR model, c: Superposition of spectral densities versus frequency, AR models.

FIGURE 3. (3.1) Nonstationary Mean, RSWOMEN Data.
Trend plus seasonal plus trading component model, a,b,c,d,e.
a: Original data and trend, b: Seasonal component, c: Trading day effect, d: Seasonal plus trading day effect, e: True, predicted and plus and minus one sigma plus predicted.
Trend plus seasonal plus AR plus trading component model, f,g,h,i,j. f: Original data and trend, g: AR component, h: Original plus trend plus AR component, i: Residual noise, j: True, predicted and plus and minus one sigma plus predicted.

FIGURE 4. Nonstationary Covariance, Seismic Data.
a. Seismic data, $y_1,\ldots,y_N$. "ordinary model" b,d,f; "intervention model" c,e,g.
b,c: $\log((y_{2m}^2 + y_{2m-1}^2)/2)$, $m=1,\ldots,N/2$ data and smoothed envelope, d,e: instantaneous power spectral density, f,g: parcors.

FIGURE 5. State Space Model Non Gaussian Discontinuous Trend Example.
a: Abruptly changing trend data, b: Smoothed state estimate, Gaussian system noise model, c: Smoothed state estimate, Pearson system-system noise model, d: Smoothed state estimate, Gaussian mixture system noise smodel, e: Posterior mean, $\pm 1,2,3\sigma$, Gaussian system noise model, f: Posterior median and (0.15,2.27,15.87,84.13,97.73) percentage points, Pearson system-system noise
model. \( g: \) Posterior median and \( (0.13, 2.27, 15.87, 84.13, 97.73) \) percentage points, Gaussian system noise model.

FIGURE 6. State Space Model Non Gaussian Envelope of Seismic Signal Example.

a: Seismic data, \( y_1, \ldots, y_N \). b: \( \log\left(\frac{y^2_m + y^2_{m-1}}{2}\right) \), \( m=1, \ldots, N/2 \) "envelope" data, c: Posterior mean, \( \pm 1, 2, 3\sigma \), Gaussian disturbances model, d: Posterior mode, \( (0.13, 2.27, 15.87, 84.13, 97.73) \) percentage points, non-Gaussian disturbances model.
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20. ABSTRACT

A variety of time series signal extraction/smoothing problems are considered from a Bayesian "smoothness priors" point of view. The origin of the subject is a smoothing problem posed by Whittaker (1923). Using a stochastic regression-linear model-Gaussian disturbances framework, we model stationary time series and nonstationary mean and nonstationary covariance time series. Smoothness priors distributions on the model parameters are expressed either in terms of time domain stochastic difference equation or frequency domain constraints. A small number of (hyper)parameters specify very complex time series behavior. The critical computation is the likelihood of the Bayesian model. Finally we show a smoothness priors state space - not necessarily Gaussian - not necessarily linear model of nonstationary time series.