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There are many areas in applied science where the non-paramotric

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THE ANALYSIS OF SOME PENALIZED LIKELIHOOD SCHEMES

by

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ABSTRACT

There are many areas in applied science where the non-parametric estimation of regression functions is important.

In this thesis a general penalized likelihood method for non-parametrically estimating regression functions under a variety of observational models is developed. The existence and numerical approximation of the estimators is studied and a cross-validatory method for estimating the smoothing parameter is presented. Implementation of the method is algorithmically straight-forward.

The procedures developed are applied to the estimation of atmospheric temperature profiles from satellite radiance data and are found to compare favorably with the currently used methodology.

A thesis under the supervision of Grace Wahba.

TABLE OF CONTENTS

1.	INTRODUCTION
	1.1 The Basic Principle
	1.2 Some Examples
	1.3 Contents of Thesis 6
	APPENDIX A. BACKGROUND MATHEMATICS 8
2.	EXISTENCE THEORY
	2.1 Introduction
	2.2 Some General Considerations
	2.3 Applications
3.	NUMERICAL ALGORITHMS
	3.1 Introduction
	3.2 Convergence Properties of the Algorithms 38
	3.3 Specific Implementations of the Algorithm51
4.	ESTIMATION OF THE SMOOTHING PARAMETER
	4.1 Introduction
	4.2 The Generalized Cross Validation Procedure57
	4.3 Some Justification for Using the GCV Estimator
5.	A BAYESIAN METHOD FOR INVERTING SATELLITE RADIANCE DATE
	5.1 Introduction
	5.2 The Bayesian Inversion Technique
	5.3 Computational Methods

5.4	Perf	formance of the Technique	•	•	•	•	.82
APPENDIX	в:	Results of Individual Retrievals	•	•	•	•	.93
APPENDIX	c:	EXISTENCE OF SMOOTH SOLUTIONS TO TH RADIATIVE TRANSFER EQUATIONS	e •	•	•	•	109
BTBLTOGR	APHY						116

a.

CHAPTER 1

INTRODUCTION

1.1 THE BASIC PRINCIPLE

The method of *Penalized Likelihood* estimation was introduced by Good and Gaskins in 1971. They argued, from a Bayesian perspective, that in order to estimate a probability density function given a random sample, $a \leq X_1, X_2, \dots, X_n \leq b$, the *penalized log likelihood functional*

$$\omega(f) = 2 \sum_{i=1}^{n} \log f(X_i) - \lambda J(f) \quad \lambda > 0$$

should be considered. J(f) is a "flamboyancy" or penalty functional such as $\int_a^b [f"(t)]^2 dt$, designed to represent prior notions about the behaviour of the density f. The estimation of the density is done by maximizing ω over a suitably chosen space of alternatives. The parameter λ controls the amount by which the data are smoothed to give the estimate.

The Good and Gaskins idea is not new to statistics, nor is it restricted to density estimation (see Kimeldorf and Wahba 1970, Leonard 1978, Silverman 1979, Tapia and Thompson 1978, Wegman 1982 and, indeed, Whittaker 1923). In this thesis some generalized versions of the penalized likelihood estimation are studied. The following specification of the method will be used: x is some function/

1

vector of parameters of interest. Conditional upon a data vector z, the estimation of x is done by minimizing the quantity:

$$I_{\lambda}(x) = Q(x | z) + \lambda J(x) \qquad \lambda > 0$$

over some set C of plausible alternatives. Q(x|z) could represent the -2 log [sampling distribution of z|x] but other choices are feasible. We give a few concrete examples.

1.2 SOME EXAMPLES

(i) Density Estimation

In 1982 Silverman introduced an interesting penalized likelihood estimator of a log probability density function. There are many density estimators which are similar in spirit (see Cox 1982 and Leonard 1978).

Given observations $a \le X_1 \le X_2 \le \dots \le X_n \le b$ Silverman's penalized likelihood estimate for the log density of the X_i 's is prescribed to be the minimizer of

$$A_0(g) = -n^{-1} \sum_{i=1}^n g(X_i) + \lambda \int_a^b [g^{(m)}(t)]^2 dt$$

over $K_1 = \{g \in W_2^m[a,b]: \int_a^b e^{g(t)} dt = 1\} m \ge 1$. It can be shown (see Silverman 1982) that minimizing A_0 over K_1 is equivalent to minimizing

$$I_{\lambda}(g) = -n^{-1} \sum_{i=1}^{n} g(X_{i}) + \int_{a}^{b} e^{g(t)} dt + \lambda \int_{a}^{b} [g^{(m)}(t)]^{2} dt$$

over $K = \{g \in W_2^m[a,b]: \int_a^b e^{g(t)} dt < \infty\}$.

(ii) Generalized Linear Regression Models

Here we have independent observations y_1, y_2, \dots, y_n where each observation y_i has a one-parameter exponential density, p, of the form:

$$p(\mathbf{y}_{i}) = \exp\{[\mathbf{y}_{i}\boldsymbol{\theta}_{i} - \mathbf{b}(\boldsymbol{\theta}_{i})]/\mathbf{a}_{i}(\boldsymbol{\phi}) + \mathbf{c}(\mathbf{y}_{i},\boldsymbol{\phi})\}$$
(*)

for suitable choices of a_i , b and c (note that ϕ , termed the scale parameter, is constant for all i). The mean and the variance of y_i can be expressed in terms of θ_i and ϕ as

$$E(y_i) = \mu_i = b'(\theta_i)$$
 and $var(y_i) = b''(\theta_i)a_i(\phi)$

The functions $a_i(\phi)$ are usually of the form ϕ/w_i (where the w_i are known). Specification of the model is completed by supplying a linearizing transformation, $g(\mu_i)$, of the mean.

$$g(\mu_i) = L_i x = \eta_i$$

The function g, referred to as the *link function*, is assumed to be monotonic and differentiable. The L_i 's are linear functionals and x is the unknown function/vector of parameters to be estimated. Penalized Likelihood estimates of x are obtained by considering I_{λ} given by:

$$I_{\lambda}(\mathbf{x}) = \sum_{i=1}^{n} [b(\theta_{i}) - y_{i}\theta_{i}]/a_{i}(\phi) + \lambda J(\mathbf{x})$$

over an appropriate space of functions.

Notice that a wide variety of distributions, including the normal, Poisson, binomial and gamma, can be written in the form of equation (*). Consequently, a great deal of useful statistical modeling can be accomplished within the Generalized Linear Interactive Modeling (GLIM) framework (see Baker and Nelder 1978). Unfortunately however, good graphical procedures for exploring the data $\{y_i\}$ do not exist (see Raynor and Bates 1983) so that the model building process can be rather slow. The primary motivation for obtaining a non-parametric determination of the regression function, x, is that it can be used as an Exploratory Data Analysis tool by the GLIM modeler.

(iii) Cox's Proportional Hazard Regression Model

Following Miller (1981), let $T_1, T_2, \ldots, T_n; C_1, C_2, \ldots, C_n$ be independent random variables. C_i is the censoring time associated with the survival time T_i . We observe $(Y_1, \delta_1), (Y_2, \delta_2), \ldots, (Y_n, \delta_n)$ where

$$Y_{i} = \min(T_{i}, C_{i}), \quad \delta_{i} = I\{T_{i} \leq C_{i}\}.$$

Covariates x_i corresponding to the survival times T_i are also available. In Cox's Proportional Hazards Model the hazard rate $\mu(t;x)$ for the distribution of T given x

4

is written as

$$\mu(t;x) = \mu_0(t)e^{\theta(x)}$$

where θ is some regression function on the covariates. Let the ordered observed times be written as

$$Y_{(1)}$$
, $Y_{(2)}$, \cdots , $Y_{(n)}$,

and $R_{(i)}$ be the risk set, i.e. the set of individuals surviving up to $y_{(i)}$. The Penalized (Partial) Likelihood estimator of θ is the minimizer of the functional

$$I_{\lambda}(\theta) = \sum_{i \in U} \{ \log \left[\sum_{j \in R_{(i)}} e^{\theta(x_j)} \right] - \theta(x_i) \} + \lambda J(\theta)$$

where U is the set of uncensored observations. The practical use of this model is currently being evaluated by Crowley and O'Sullivan (1983), a report should be forthcoming.

(iv) Normal Regression Model

The observed data $\{z_i\}_{i=1}^n$ are functionals of x contaminated by white noise.

$$z_i = n_i(x) + \varepsilon_i \quad \varepsilon \sim N(0,\sigma)$$

where n_i 's are functionals of x. x is estimated by minimizing:

$$I_{\lambda}(x) = \frac{1}{n} \sum_{i=1}^{n} [z_{i} - \eta_{i}(x)]^{2} + \lambda J(x), \quad \lambda > 0.$$

Apart from the data smoothing applications of this model (see Kimeldorf and Wahba 1970), there are a number of inversion problems in engineering where the normal regression assumption becomes a useful approximation (see Westwater 1979).

1.3 CONTENTS OF THESIS

There are a number of questions which arise in connection with the actual implementation of the penalized likelihood method. The penalized likelihood estimate is the minimizer of the penalized likelihood functional, I_{λ} , over an appropriate space of functions. The technical issues addressed in the following chapters include:

- Existence of the penalized likelihood estimate.

- Numerical computation of the estimates for fixed λ .

- Choice of the smoothing parameter $\boldsymbol{\lambda}$.

The analysis is carriedout first for the general case and then applied to various particular models. Existence questions are dealt with in Chapter 2; the approach is similar to that taken by Wegman (1982). Numerical methods are presented in Chapter 3 - versions of the Newton-Raphson and Gauss-Newton algorithms are analyzed. The material in these latter two chapters involves a fair bit of mathematical machinery, the elements of which are contained in an appendix at the end of this chapter. Generalized Cross Validation-type estimators of the smoothing parameter λ are proposed and partially justified in Chapter 4, while the final chapter features an application of the methodology to the remote sensing of atmospheric temperature profiles. An asymptotic analysis of penalized likelihood estimators is currently under study (see Cox and O'Sullivan 1983).

7

APPENDIX A. BACKGROUND MATHEMATICS

NOTATION

H is a real Hilbert space with inner product $\langle \cdot, \cdot \rangle$ and norm $\|\cdot\|$ (so $\|x\|^2 = \langle x, x \rangle |x | \epsilon$ H). The dual space of continuous linear functionals on H is denoted by H*. $\|\cdot\|$ induces a norm $\|\cdot\|_{*}$ on H* given by

$$\| \mathfrak{L} \|_{\star} = \sup_{\mathbf{x} \in \mathbf{H}} \frac{| \mathfrak{L}(\mathbf{x}) |}{\| \mathbf{x} \|} \ \mathfrak{L} \in \mathbf{H}^{\star}.$$

For every $l \in H^*$ $\exists y_l \in H$ such that

 $\ell(x) = \langle y_{\ell}, x \rangle \quad \forall x \in H$

The relation $\ell \iff y_{\ell}$ establishes an isometric isomorphism between the spaces H and H^{*}. We also have that H^{*} is a Hilbert space with inner product:

$$\langle \ell, k \rangle_{\star} = \langle y_{0}, y_{k} \rangle \qquad \ell, k \in H^{\star}$$

Finally if H_1 and H_2 are two Hilbert spaces then the space of continuous linear operators from H_1 into H_2 will be denoted by $L(H_1, H_2)$.

THE WEAK TOPOLOGY

In the usual or norm topology on H, an open ball, B_r , of radius r about the origin is defined by:

 $B_r = \{x \in H : ||x|| < r\}$

However the space H* generates another topology on H called the *weak* topology: this is the weakest topology on H with respect to which the elements of H* are continuous on H.

The usual topological notions of convergence, compactness, closedness etc. will be used in many places (see Simmons 1963). Whenever we are talking about these notions, unless explicit reference is made to the weak topology, the norm topology will be what we have in mind. Thus when a sequence $\{x_n\}$ converges to a point x in H we will write:

$$x_n \rightarrow x$$
 as $n \rightarrow \infty$

if norm convergence is intended $(\lim_{n \to \infty} ||x_n - x|| = 0)$. Whereas if $\{x_n\}$ converges weakly to x we will write:

$$x \rightarrow x as n \rightarrow \infty$$

meaning $\lim_{n \to \infty} |\ell(\mathbf{x}_n) - \ell(\mathbf{x})| \to 0 \quad \forall \ \ell \in \mathrm{H}^*$.

PROPERTIES OF NON-LINEAR FUNCTIONALS ON H

We will be interested in various properties of nonlinear functionals defined on H, the most important of these being convexity, continuity and differentiability.

9

CONVEXITY

Definition 1.A.1. A set C in H is convex if ∀ x, y ∈ C

$$\lambda x + (1 - \lambda) y \in C$$
 $0 < \lambda < 1$.

Definition 1.A.2. A functional f on H is convex if $\forall x, y \in H$ and $0 < \lambda < 1$

$$f(\lambda x + (1 - \lambda)y) \leq \lambda f(x) + (1 - \lambda)f(y).$$

f is strictly convex if the above inequality is strict whenever $x \neq y$.

Definition 1.A.3. The epigraph of a functional f on H is the set:

epi f = { (u,a) \in H × R: f(u) \leq a }.

It is the set of points in $H \times R$ which lie above the graph of f. The notions of convex sets and convex functions are intimately related; a function is convex if and only if it's epigraph is convex. One important property of convex functions on finite-dimensional spaces is the following:

Lemma 1.A.1. Suppose H is finite dimensional. If f is strictly convex and bounded below and has a minimum in H then $f(x) \rightarrow \infty$ as $||x|| \rightarrow \infty$ (i.e. f is coercive on H). Proof: See Rockafellar (1970). Definition 1.A.4. A function f is quasi-convex if $\forall x, y \in H$ and $0 < \lambda < 1$

 $f(\lambda x + (1 - \lambda)y) \le max (f(x), f(y))$

f is strongly quasi-convex if the above inequality is strict whenever $x \neq y$.

It is easy to see that if f is convex or quasi-convec then a unique local minimizer of f is the unique global minimizer of f. Also if f is strictly convex or strongly quasi-convex any minimizer of f is the unique global minimizer.

Definition 1.A.5. A function f is uniformly convex on H if $\forall x \neq y$

 $f(\frac{x+y}{2}) \leq \frac{1}{2}f(x) + \frac{1}{2}f(y) - \delta(||x-y||)$

for some real-valued, continuous, monotone function δ ; $\delta(t) \ge 0$ for $t \ge 0$ with $\delta(t) = 0$ if and only if t = 0.

Local versions of convexity, quasi-convexity, uniform convexity etc. can also be defined. For instance f is locally convex at $x \in H$ if f is convex in some neighborhood of x. We get the following lemma.

Lemma 1.A.2. If f is a convex function and x^* is a minimizer of f then if f is locally uniformly convex (or locally strictly convex) at x^* then x^* is the unique

global minimizer of f.

CONTINUITY AND LOWER SEMI-CONTINUITY

Definition 1.A.6. A function f is (weakly) continuous at \overline{u} if for all sequences, $\{u_n\}$, (weakly) converging to \overline{u} , $\lim_{n \to \infty} f(u_n) = f(\overline{u})$.

Definition 1.A.7. A function f is (weakly) *lower semi*continuous on H if it satisfies one of the two equivalent conditions:

 $\forall a \in R \{ u \in H : f(u) \leq a \}$ is (weakly) closed.

or

 $\overline{u} \in H$, liminf $f(u_n) \ge f(\overline{u})$

for all sequences (weakly) converging to \bar{u} . A further characterization of lower semi-continuity is also available.

Theorem 1.A.3. A function $f: H \rightarrow R$ is (weakly) lower semi-continuous if and only if it's epigraph is (weakly) closed.

The concept of lower semi-continuity plays an important role in minimization problems because of the following result. <u>Theorem 1.A.4</u>. Let f be a weakly lower semi-continuous functional defined on a weakly compact set C; then there exists x_0 in C such that $f(x_0) = \inf_{x \in C} f(x)$.

Proof:

Let $m = \inf_{x \in C} f(x)$; then there exists a sequence $\{x_n\}$ in C such that $\lim_{n \to \infty} f(x_n) = m$. Since C is weakly compact, there is a weakly convergent subsequence $\{x_{n_i}\}$ of $\{x_n\}: x_{n_i} \xrightarrow{W} x_0$ for some x_0 in C. But by weak lower semi-continuity and the fact that $\lim_{n_i \to \infty} f(x_{n_i}) = m$ we have

$$f(x_0) \leq \liminf_{n_i \to \infty} f(x_{n_i}) = m \leq f(x_0)$$

Thus $f(x_0) = m$.

The above theorem has two important corollaries.

Corollary 1.A.5. A weakly lower semi-continuous functional achieves its infimum on every closed and bounded convex subset of H.

Proof:

In a real Hilbert space, every closed convex set is weakly closed (see Ekeland and Teman (1973) pp. 3-7) while a weakly closed and bounded set is weakly compact. Therefore Theorem 1.A.4 applies.

Corollary 1.A.6. Let f be a weakly lower semi-continuous functional on an unbounded closed convex set C. Suppose lim $f(x) = \infty$ in C. Then f achieves its in- $||x|| \rightarrow \infty$ fimum on C.

Proof:

If $f(x) = \infty \forall x \in C$ then we are done. Otherwise for M large enough inf $f = \inf_C f$, where B is the unit ball. The result now follows by corollary 1.A.5 since CoMB is a closed and bounded convex set.

FRECHET DIFFERENTIABILITY

Let f be a mapping from a Hilbert space H_1 (norm $\|\cdot\|_1$) into a Hilbert space H_2 (norm $\|\cdot\|_2$).

Definition 1.A.8. f is Frechet differentiable at $x \in H_1$ if there exists a bounded linear operator $f'(x) \in L(H_1, H_2)$ satisfying:

$$\lim_{\|y\|_{2} \to 0} \frac{\|f(x+y) - f(x) - f'(x)y\|_{1}}{\|y\|_{2}} = 0$$

f'(x) is called the Frechet derivative of f at x. Higher order Frechet derivatives are analogously defined; the second Frechet derivative of f at x, f"(x), is an element of $L(H_1, L(H_1, H_2))$ etc.

Versions of the usual mean value theorems apply to Frechet differentiable mappings. A couple of these results, which we will use later on, are worth noting. <u>Theorem 1.A.7</u>. Let $f : D \subset H \rightarrow R$, and assume that f has a second Frechet derivative at each point of a convex set $D_0 \subset D$. Then, for any $x, y \in D_0$, there is a $t \in [0,1]$ such that

$$f(y) - f(x) - f'(x) (y-x) = \frac{1}{2}f''(x+t(y-x))(y-x)(y-x)$$

Theorem 1.A.8. Assume that $f : D \subset H \rightarrow R$, has a second Frechet derivative at $x \in D$. Then

 $\lim_{\|h\| \to 0} (1/\|h\|^2) [f(x+h) - f(x) - f'(x)h - \frac{1}{2}f''(x)hh] = 0$

for any $h \in H$.

Definition 1.A.9. A function f on H is uniformly positive definite if for some positive constant M we have:

 $f''(x)hh \ge M \|h\|^2$

for all $x, h \in H$.

Functions which are uniformly positive definite on H clearly satisfy: $f(x) \rightarrow \infty$ as $||x|| \rightarrow \infty$.

THE GRADIENT AND HESSIAN OF A FUNCTIONAL

The usual notions of gradients and hessians of functionals on R^n generalize to real Hilbert spaces. Suppose g: H \rightarrow R is a functional on H. The first and second Frechet derivatives of g at a point x in H are denoted by g'(x) and g''(x) respectively. Now since $g'(x) \in H^*$ and $g''(x) \in L(H,H^*)$, by the isometric isomorphism between H and H^{*}, $\exists \cdot \nabla g(x) \in H$ and $H_g(x) \in L(H,H)$ with the properties that:

 $\langle \nabla g(\mathbf{x}), \mathbf{y} \rangle = g'(\mathbf{x}) \mathbf{y} \forall \mathbf{y} \in \mathbf{H}$

 $\|\nabla g(x)\| = \|g'(x)\|_{*}$

and

$$< H_{g}(x)y, z > = g''(x)yz \forall y, z \in H$$

 $\|H_{g}(x)y\| = \|g'(x)y\|_{*} \forall y \in H$

In $\mathbb{R}^n \nabla g(x)$ and $\mathbb{H}_g(x)$ correspond to the gradient and hessian, respectively.

CHAPTER 2

EXISTENCE THEORY

2.1 INTRODUCTION

In this chapter we give conditions under which the Penalized Likelihood estimators will exist. Attention will be restricted to the case where the penalty functional J is, at least, convex (indeed, more often than not, J(x) = $||Px||^2$ where P is a projection operator with finite dimensional null space). One of the main results, Theorem 2.2.5 and its corollary, gives conditions when the existence of a minimizer of Q(x|z) guarantees the existence of the Penalized Likelihood estimator; this related to Silverman's 1982 Annals of Statistics conjecture. We begin with some general considerations and go on to discuss some applications in section 3. The major references for this chapter are Daniel (1971), Ekeland and Teman (1973) and Ortega and Rheinbold (1970). The notation introduced in Chapter 1 will be employed throughout.

2.2 SOME GENERAL CONSIDERATIONS

Consider the functional:

 $I_{\lambda}(x) = Q(x) + \lambda J(x), \quad \lambda > 0.$

17

Our program is to find sufficient conditions for there to exist minimizers of I_{λ} in closed convex unbounded subsets C of H. The next theorem, a restatement of corollary 1.A.6., clarifies this problem somewhat.

Theorem 2.2.1.

Let C be a closed convex subset of a Hilbert space H. Suppose $I_{\lambda}: C \rightarrow R$ is coercive (i.e. $I_{\lambda}(x) \rightarrow \infty$ as $||x|| \rightarrow \infty, x \in C$) and is weakly lower semi-continuous on C, then I_{λ} attains its infimum on C.

Proof: [See corollary 1.A.6]

In the context of this theorem we now ask ourselves what conditions must the functionals Q and J satisfy in order that

(i) I_{λ} is weakly lower semi-continuous on C?

(ii) I $_\lambda$ is coercive on C?

The first of these can be answered as follows.

Theorem 2.2.2.

Suppose C is a closed convex subset of H. Let Q be weakly continuous on C, and J (bounded below) continuous and convex on C. Then I $_{\lambda}$ is weakly lower semicontinuous.

Proof:

J is continuous and convex so it is weakly lower continuous (see the appendix of Chapter 1). Now since $\lambda > 0$ and Q is weakly lower semi-continuous (it is weakly continuous) it follows that $Q + \lambda J$ is weakly lower semicontinuous.

Conditions guaranteeing the coerciveness of I_{λ} on general convex sets C are not as easy to identify. However, when Q and J are well behaved, and C satisfies some mild conditions, the problem becomes more tractable. Suppose H can be split into two subspaces, H_0 and H_1 , with $H = H_0 \oplus H_1 : H_0$ is finite-dimensional. Let $J(x) = ||Px||^2$ where P is a projection operator on H with null space H_0 and range H_1 . We consider subsets C of H which satisfy the following property:

Definition 2.2.3.

A convex set C in $H_0 \oplus H_1$ satisfies property 1 if for any sequence in C for which $\|(I-P)x_n\| \to \infty$ and $\{\|Px_n\|\}$ is bounded \exists sequences $\{\mu_n\} \subset [0,1], \{u_n\} \subset H_0, \{y_n\} \subset C$ and $z \in H$ for which

$$u_n + z = \mu_n x_n + (1 - \mu_n) y_n$$

with $\|u_n\| \to \infty$ and $\{y_n\}$ bounded.

Definition 2.2.4.

A convex set C in $H_0 \oplus H_1$ satisfies property 2 if it satisfies property 1 and the choice of z does not depend on the sequence $\{x_n\}$. Convex sets satisfying property 1 are common in applications; for instance in $W_2^m[a,b]$ with $H_0 = \{x \in W_2^m[a,b]: x^{(m)} \equiv 0\}$ the sets $C_{q,\alpha} = \{x \in W_2^m[a,b]: x^{(q)} \ge \alpha\}, 0 \le q < m$ and α some continuous function. To see this let $\{x_n\} \in C_{q,\alpha}$ be such that $\|Px_n\|^2 = \int_a^b [x_n^{(m)}(t)]^2 dt$ is bounded and $\|(I-P)x_n\| \equiv \|P_0x_n\|^2 \to \infty$. Since $\|Px_n\|$ is bounded, q < m and α is continuous, Ξ a polynomial, M(t), such that $-M^{(q)}(t) < (Px_n)^{(q)}(t) < M^{(q)}(t) t \in [a,b],$ and $\alpha < M$. Let $u_n = \frac{1}{2} P_0 x_n, y_n = 2M - Px_n$ and z = 2M. By definition, $\{y_n\}$ is a bounded sequence in $C_{q,\alpha}$,

and since

$$u_n + z = \frac{1}{2}x_n + \frac{1}{2}y_n$$

we have that $C_{q,\alpha}$ indeed satisfies property 1.

Theorem 2.2.5.

Let C be a convex subset of H satisfying property 1. Suppose Q satisfies the following conditions:

(i) coercive on $[x + H_0] \cap C$, for $x \in H$

(ii) quasi-convex on C

(iii) bounded on bounded subsets of C

(iv) bounded below on C

then I_1 is coercive on C.

Proof:

$$I_{\lambda}(\mathbf{x}) = Q(\mathbf{x}) + \lambda \|P\mathbf{x}\|^{2}$$

If I_{λ} is not coercive on C then there is a sequence $\{x_n\}$ in C for which $\{I_{\lambda}(x_n)\}$ is bounded but $\{x_n\} \to \infty$ as $n \to \infty$.

Case 1: $\{x_n\}$ has a subsequence, $\{x_{n_k}\}$, for which $\|Px_{n_k}\| \to \infty$.

Here, since Q is bounded below on C, we have that $I_{\lambda}(x_{n_{k}}) = Q(x_{n_{k}}) + \lambda \|Px_{n_{k}}\|^{2} \to \infty, \text{ contradicting the definition of } \{x_{n}\}.$

Case 2: $\{\|Px_n\|\}$ is bounded.

C satisfies property 1, therefore \exists sequences $\{\mu_n\} \subset [0,1], \{u_n\} \subset H_0, \{y_n\} \subset C$ and $z \in H$ for which

$$u_n + z = \mu_n x_n + (1 - \mu_n) y_n$$

with $\|u_n\| \to \infty$ and $\{y_n\}$ bounded. Since Q is quasiconvex we have that:

$$Q(u_n + z) \leq max(Q(x_n), Q(y_n))$$

Since Q is coercive on $[z + H_0] \cap C$, and bounded on bounded subsets of C it follows that $Q(x_n) \rightarrow \infty$ which of course contradicts the definition of $\{x_n\}$.

Corollary 2.2.6.

Let C be a convex set satisfying property 1 and suppose Q is bounded below and bounded on bounded subsets of C. If Q is strictly convex on C and has minimizers on subsets of C of the form $[x + H_0] \cap C \ x \in H$, then I_1 is coercive on C.

Proof:

Since Q is strictly convex and has a minimum on $[x + H_0] \cap C$ Q is coercive on sets of this form (see Lemma 1.A.1).

The above results obviously hold when the set C satisfies property 2: the coerciveness of Q need then only be checked on $[z + H_0] \cap C$ for a fixed z in H.

2.3 APPLICATIONS

Consider the following set up $H = H_0 \oplus H_1$ with H_0 finite-dimensional, P is a projection operator on H with range space H_1 and null space H_0 . $J(x) = ||Px||^2$, and L_i/N_i are continuous linear/non-linear functionals on H. Finally, C is a closed convex subset of H satisfying property 1 of section 2. Within this framework we investigate the existence of minimizers of the I_{λ} 's arising in a few specific models.

(i) SILVERMAN'S PENALIZED LIKELIHOOD DENSITY ESTIMATES

Given observations $a \le X_1 \le X_2 \le \cdots \le X_n \le b$ Silverman's (1982) penalized likelihood estimate for the log density of the X_i 's is prescribed to be the minimizer of

$$A_0(g) = -n^{-1} \sum_{i=1}^{n} g(X_i) + \lambda \int_a^b [g^{(m)}(t)]^2 dt$$

over $C_1 = \{g \in W_2^m[a,b]: \int_a^b e^{g(t)} dt = 1\} m \ge 1$. It can be shown (see Silverman 1982) that minimizing A_0 over C_1 is equivalent to minimizing

$$A(g) = -n^{-1} \sum_{i=1}^{n} g(X_i) + \int_a^b e^{g(t)} dt + \lambda \int_a^b [g^{(m)}(t)]^2 dt$$

over $C = \{g \in W_2^m[a,b]: \int_a^b e^{g(t)} dt < \infty\}.$

Since $m \ge 1$, if $g \in W_2^m[a,b]$ then |g| is bounded and so $\int_a^b e^{g(t)} dt$ is finite. Consequently $C = W_2^m[a,b]$. It follows trivally that C satisfies property 2 of section 2 with $z \equiv 0$.

A is strictly convex and bounded below by zero on C. Since evaluation is continuous on W_2^m for $m \ge 1$ and since $\int_a^b e^g dt$ is weakly continuous (by continuity of e^x), $-n^{-1} \sum_{i=1}^n g(X_i) + \int_a^b e^{g(t)} dt$ is weakly continuous and consequently bounded on bounded subsets of C. Therefore, by Theorem 2.2.2, A is weakly lower semi-continuous and so by the Corollary 2.2.6 A has a minimizer in C provided $-n^{-1} \sum_{i=1}^n g(X_i) + \int_a^b e^{g(t)} dt$ has a minimizer in {1,t,t²,...,t^{m-1}}. This provides an alternative justification of Theorem 4.1 in Silverman's (1982) paper.

(ii) GENERALIZED LINEAR MODELS

Wedderburn (1976) discusses the existence and uniqueness of the maximum likelihood estimates for certain generalized linear models over finite-dimensional parameter spaces. Here we use some of his results to discuss two particular models in some detail.

Log-linear Poisson regression

This model is especially useful for studying systems in which the measurement process is subject to shot noise (see Nicot and Yandell (1983) for an interesting application involving the mapping of fungal activity in potato fields). Barndorff-Nielsen (1978) discusses a parametric regression form of the model (see also Kalbfleisch and Sprott 1974). The data consist of independent counts $\{y_i, i = 1, 2, ..., n\}$. Each $y_i \sim P(\mu_i)$ and μ_i is modeled log-linearly as:

 $\mu_{i} = e^{L_{i}x}$ i = 1, 2, ..., n

where L_i 's are continuous linear functionals of x. Thus the likelihood of x is given by:

$$\Phi(\mathbf{x}) = \prod_{i=1}^{n} \frac{\mu_{i}^{y_{i}} e^{-\mu_{i}}}{Y_{i}!}$$

from which it follows that I_{λ} has the form:

$$I_{\lambda}(\mathbf{x}) = \sum_{i=1}^{n} \{ e^{\mathbf{L}_{i}\mathbf{x}} - y_{i}\mathbf{L}_{i}\mathbf{x} + \log(y_{i}!) \} + \lambda \|\mathbf{P}\mathbf{x}\|^{2}$$

Let dim(H₀) = m; then $\exists \{\phi_j\}_{j=1}^m : H_0 = \text{span } \{\phi_1, \phi_2, \dots, \phi_m\}$. Now given $x \in [z + H_0]$

$$x = \sum_{j=1}^{m} \beta_{j} \phi_{j} + z$$

Let $B = \{\beta \in R^m: \sum_{j=1}^m \beta_j \phi_j + z \in [z + H_0] \cap C\}$. B is a closed convex subset of R^m . Let T be a matrix with ijth element $L_i \phi_j$ i = 1,2,...,n j = 1,2,...,m, then

 $L_i x = T_i^{\dagger} \beta + L_i z$

fin :

where $T_i = (L_i \phi_1, L_i \phi_2, \dots, L_i \phi_m)'$. If rank $(T) \ge m$ then it follows easily that I_λ is strictly convex on H. Also, since the L_i 's are continuous linear functionals, I_λ is bounded on bounded subsets of H and so, by Theorem 2.2.2., I_λ is weakly lower semi-continuous $(||Px||^2)$ is continuous and convex and $\sum_{i=1}^{n} \{e^{L_i X} - y_i L_i x + \log(Y_i!)\}$ is weakly continuous and bounded below). Consequently, by the corollary to Theorem 2.2.5, I_λ will achieve its minimum on C if, dropping extraneous terms, $\sum_{i=1}^{n} \{e^{L_i X} - y_i L_i X\}$ has a minimum i = 1 $\sum_{i=1}^{n} \{e^{L_i X} - y_i L_i X\}$ has a minimum i = 1 $\sum_{i=1}^{n} \{e^{L_i X} - y_i L_i X\}$ has a minimum i = 1 $\sum_{i=1}^{n} \{e^{L_i X} - y_i L_i X\}$ has a minimum i = 1 $\sum_{i=1}^{n} \{e^{T_i \beta} - y_i^* T_i' \beta\}$ has a minimum in B, where $y_i^* = y_i / e^{L_i Z}$. Therefore we need only give conditions under which Q has a minimum in B.

Let T^0 be the matrix obtained from T by removing the rows corresponding to zero y_i counts. Following Wedderburn (1976) we can deduce:

Lemma 2.3.1.

If rank (T) \ge m and the row space of T⁰ is the same as the row space of T, then Q(β) has a minimum in B.

Proof:

Q is convex and continuous on B, so, we are done once we establish that Q is coercive on B. Suppose $\|\beta\| \to \infty$ in B. Then since T is of full rank $\|T\beta\| \to \infty$, but this in turn implies that $\|T^{0'}\beta\| \to \infty$, that is $\sum_{\substack{i:y,\neq 0\\i:y,\neq 0}} \{e^{T_{i}\beta} - y_{i}^{*}T_{i}^{*}\beta\} \to \infty$. It follows, arguing by coni:y_i \neq 0

Therefore, under the hypotheses of the above lemma, I_{λ} achieves its minimum in C. The uniqueness of this minimum follows by the strict convexity of I_{λ} on C.

The logistic regression model

Cox (1970) discussed this model as a special case of more general quantal response models. The data in these models (in the form of y_i individuals responding out of n_i) are considered as independent binomial variates whose

means are functionally related via a common regression.

In the logistic case we observe $\{y_i, i = 1, 2, ..., n\}$ where $y_i \sim B(n_i, p_i)$ and the logits are written as:

$$\log \left[\frac{p_{i}}{1-p_{i}}\right] = L_{i} \times i = 1, 2, ..., n$$

Silverman (1978), Raynor and Bates (1983) and Villalobos (1983) have considered this model useful in the analysis of discrimination problems. The likelihood of x given the observed data is:

$$\Phi(\mathbf{x}) = \prod_{i=1}^{n} p_{i}^{\text{Yi}} (1 - p_{i})^{n_{i} - \text{Yi}}$$

so that I, has the form:

$$I_{\lambda}(\mathbf{x}) = \sum_{i=1}^{n} [n_{i} \log (1 + e^{L_{i}\mathbf{x}}) - y_{i}L_{i}\mathbf{x}] + \lambda \|P\mathbf{x}\|^{2}$$

Differentiating w.r.t. x we get for any $h \in H$:

$$I_{\lambda}''(x)hh = \sum_{i=1}^{n} n_{i} \frac{e^{L_{i}x}}{(1+e^{L_{i}x})^{2}} (L_{i}h)^{2} + 2\lambda \|Ph\|^{2}$$
(*)

= 0 iff
$$\sum_{i=1}^{11} (L_i h)^2 = 0$$
 and Ph = 0.

Let T be defined as in the Poisson model and suppose we are interested in minimizing I_{λ} over the whole of H. Equations (*) imply that I_{λ} is strictly convex provided rank (T) \geq m. I_{λ} is obviously bounded below, the continuity of the L_i 's implies that I_{λ} is weakly lower semi-continuous and hence bounded on bounded subsets of H. It follows from Corollary 2.2.6 that when rank $(T) \ge m I_{\lambda}$ has a minimizer in H whenever $\sum_{i=1}^{n} [n_i \log (1+e^{T_i \cdot \beta}) - y_i T_i \cdot \beta]$ has a minimum in \mathbb{R}^m . If the rank $(T) \ge m$ then $\sum_{i=1}^{n} [n_i \log (1+e^{T_i \cdot \beta}) - y_i T_i \cdot \beta]$ on \mathbb{R}^m constitutes a regular exponential family distribution of order m (see Barndorff-Nielsen 1978) and using the Barndorff-Nielsen theory of such families we obtain the following result.

Lemma 2.3.2.

Let $S = \{s = (s_1, \dots, s_{\nu}, \dots, s_m) : s_{\nu} = \sum_{i=1}^{n} y_i L_i \phi_{\nu};$ $y_i = 0, 1, \dots, n_i\}$ and C_s be the closure of the convex hull of S. If rank $(T) \ge m$ then $\sum_{i=1}^{n} [n_i \log (1 + e^{T_i \beta}) - y_i T_i^{i} \beta]$ has a minimum in R^m if and only if the observed value of s lies in the interior of C_s .

Proof:

The condition on the rank of T guarantees that $\sum_{i=1}^{n} [n_i \log (1 + e^{T_i \beta}) - y_i T_i \beta] \text{ constitutes a regular exponential}$ family distribution of order m and the result follows by appealing to Corollary 0.6, p. 153 of Barndorff-Nielsen (1978). \Box

Therefore, under the hypotheses of Lemma 2.3.2, I $_{\lambda}$ has a minimizer in H. The uniqueness of this minimizer follows from the strict convexity of I $_{\lambda}$ on H. For the case in which L_i's are evaluations and the n_i's are all l's, Silverman (1978) and Villalobos (1983) have obtained a more intuitive characterization of the result in Lemma 2.3.2.

(iii) Cox's Proportional Hazard Regression Model

Suppose the explanatory variables x_1, x_2, \ldots, x_n lie in [a,b]. Following section 2 of Chapter 1, let the Penalized (Partial) Likelihood estimator of the proportional hazard regression function, θ , be the minimizer of the functional

$$\begin{split} \mathbf{I}_{\lambda}(\theta) &= \sum_{i \in \mathbf{U}} \{ \log \left[\sum_{j \in \mathbf{R}_{(i)}} e^{\theta \left(\mathbf{x}_{j} \right)} \right] - \theta \left(\mathbf{x}_{i} \right) \} + \lambda \int_{a}^{b} \left[\theta^{\left(m \right)} \left(\mathbf{x} \right) \right]^{2} d\mathbf{x} \\ \text{over } \mathbf{C} &= \{ \theta - \theta \left(a \right), \ \theta \in \mathbf{W}_{2}^{m} [a, b] \} \text{ for } \mathbf{m} \geq 1. \quad (\text{The } -\log \left\{ \text{partial likelihood} \} \text{ evaluated at } \theta - \theta \left(a \right) \text{ is equal to the} \\ - \log \left\{ \text{partial likelihood} \right\} \text{ evaluated at } \theta. \text{ Hence, the motiva-} \\ \text{tion minimizing } \mathbf{I}_{\lambda} \text{ over C.} \text{ Clearly } \mathbf{C} \text{ satisfies pro-} \\ \text{perty 2 of section 2 with } "z \equiv 0". \text{ Also } \mathbf{I}_{\lambda} \text{ is bounded} \\ \text{below by zero and since evaluation is continuous } (\mathbf{m} \geq 1) \\ \text{the } -\log \left\{ \text{partial likelihood} \right\} \text{ is weakly continuous.} \end{split}$$

With $p_{ij}(\theta) = \frac{e^{\theta(x_j)}}{\sum e^{\theta(x_k)}}$ we have for any $\theta, \phi \in C$: $k \in \mathbb{R}(i)$

$$I_{\lambda}"(\theta) \phi \phi = \sum_{i \in U} [\sum_{j \in R} p_{ji}(\theta) \phi^{2}(x_{j})]$$
$$- [\sum_{j \in R} p_{ji}(\theta) \phi(x_{j})]^{2}]$$
$$+ 2\lambda \|P\phi\|^{2} \ge 0$$

= 0 iff
$$\phi(x_j) = 0$$
 for $j \in \bigcup_{i \in U} R(i)$

and
$$P\phi = 0$$
.

It follows that I_{λ} is strictly convex on C provided there are at least m distinct x_j 's in $\bigcup_{i \in U} R_{(i)}$. Hence, by the standard arguments, I_{λ} has a minimum in C provided there exist $\beta \in \mathbb{R}^{m-1}$ to minimize the finite dimensional -log {partial likelihood}:

$$Q(\beta) = \sum_{i \in U} \{ \log \left[\sum_{j \in R_{(i)}} e^{Tj\beta} \right] - T_{i}\beta \}$$

where $T_{i} = (x_{i}, x_{i}^{2}, \dots, x_{i}^{m-1})$.

(iv) THE NON-LINEAR REGRESSION MODEL

The penalized likelihood functional, ${\rm I}_{\lambda},$ is given by

$$I_{\lambda}(\mathbf{x}) = \sum_{i=1}^{n} [z_i - N_i(\mathbf{x})]^2 + \lambda \|P\mathbf{x}\|^2$$

where N_i 's are non-linear functionals. If the N_i 's are

weakly continuous, then I $_{\lambda}$ is weakly lower semi-continuous and bounded on bounded subsets of H.

We need to work harder here to develop meaningful conditions under which I_{λ} is coercive on H. The conditions which we come up with are strongly influenced by our expierence with the radiative transfer equations (see Chapter 5 and Appendix B). We begin by considering *comparison* functionals I_{ϕ} defined on C in the following manner.

 $I_{\Phi}(x) = \Phi(N_{1}(x), N_{2}(x), \dots, N_{n}(x)) + ||Px||^{2}$ (**)

where $\Phi: \mathbb{R}^n \to \mathbb{R}$ is bounded below and satisfies $\Phi(\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n) \leq \Phi(\mathbf{x}_1^*, \mathbf{x}_2^*, \dots, \mathbf{x}_n^*)$ whenever $|\mathbf{x}_i| \leq |\mathbf{x}_i^*|$ \forall i,i = 1,2,...,n.

These comparison functionals will be used to establish the coerciveness of I_{λ} for certain classes of non-linear functionals N_i . Whenever I_{Φ} is coercive then so is the corresponding I_{λ} ; the next lemma is needed to prove this assertion.

Lemma 2.3.3.

If $\{x_k\}$ is any sequence in \mathbb{R}^n for which $\Phi(x_k) \neq \infty$ as $k \neq \infty$ then $\|x_k\|_{E_n}$ (the Euclidean norm of x_k) also tends to infinity as $k \neq \infty$.

Proof:

Suppose \exists a bounded subsequence of $\{x_k\}, \{x_k\}$ say, then by definition of Φ , $\{\Phi(x_k)\}$ must also be bounded. However, this contradicts the fact that $\Phi(x_k) \neq \infty$ as $k \neq \infty$. Therefore, $\forall M \exists k_M$ such that $\forall k \geq k_M$ $\|x_k\| \geq M$ i.e. $\lim_{k \neq \infty} \|x_k\| = \infty$. \Box

Theorem 2.3.4.

Let C be an unbounded convex set and suppose ${\rm I}_{\Phi}$ given by (**) is coercive on C, then ${\rm I}_{\lambda}$ is also coercive on C.

Proof:

Suppose not, then for some $M > 0 = \{x_{l}\} \subseteq C$ with the property that

$$I_{\lambda}(x_{0}) \leq M$$
 for $||x_{0}|| > \ell$ $\ell = 1, 2, \ldots$

For this sequence we must have, by definition of I_{λ} , that

$$\| \operatorname{Px}_{\ell} \| \leq \frac{1}{\lambda} \, \mathsf{M} \, \forall \, \ell$$

However, the coerciveness of I $_{\Phi}$ on C implies that

$$\Phi(N_1(x_{\ell}), N_2(x_{\ell}), \dots, N_2(x_{\ell})) \rightarrow \infty \text{ as } \ell \rightarrow \infty$$

Hence by Lemma 2.3.3.

$$\sum_{i=1}^{n} [N_{i}(x_{\ell})]^{2} \rightarrow \infty \Rightarrow \sum_{i=1}^{n} [z_{i} - N_{i}(x_{\ell})]^{2} \rightarrow \infty \Rightarrow I_{\lambda}(x_{\ell}) \rightarrow \infty$$
This contradicts the definition of $\{x_{\ell}\}$. Therefore $\forall M \exists \ell$ such that $\forall x \in C$ with $||x|| > \ell$ $I_{\lambda}(x) \ge M$ i.e. I_{λ} is coercive on C.

We finish this section by presenting a theorem which provides a set of conditions under which the functionals N_i yield coercive I_λ 's on C. This result will be used to study the existence of smooth solutions to the ill-posed problem associated with the radiative transfer equations (see Appendix B).

Theorem 2.3.5.

Let C be a convex subset of H satisfying property 1. Suppose $\exists \phi : R \rightarrow R$ which is monotonic increasing (and bounded below) in the modulus of its argument. Suppose further that

(s.1) $\sum_{i=1}^{n} \phi[N_{i}(x)] \text{ is strictly convex on } C$ (s.2) $\sum_{i=1}^{n} \phi[N_{i}(x)] \text{ has a proper minimum on}$ $[x + H_{0}] \cap C, x \in H.$

Then I_{λ} is coercive on C.

 $\frac{\text{Proof:}}{\text{Let } \Phi(\mathbf{x}) = \sum_{i=1}^{n} \phi(\mathbf{x}_{i}). \text{ Then } \Phi \text{ clearly satisfies}$ properties (**) above. Next define I_{Φ} as follows:

33

$$I_{\Phi}(x) = \sum_{i=1}^{n} \phi[N_{i}(x)] + ||Px||^{2}$$

By (s.1) I_{Φ} is convex on C and since the N_i 's are weakly continuous, I_{Φ} is bounded on bounded subsets of H. Finally, since $\sum_{i=1}^{n} \phi[N_i(x)]$ is convex and has a i=1 $(x + H_0] \cap C \times \in H$, it follows that I_{Φ} , is coercive on C and so by Theorem 2.3.4. that I_{λ} is coercive on C. \Box

CHAPTER 3

NUMERICAL ALGORITHMS

3.1 INTRODUCTION

This chapter discusses some numerical methods for minimizing the penalized likelihood functionals. Recall that penalized likelihood functionals have the form:

 $I_{\lambda}(x) = Q(x|z) + \lambda J(x), \lambda > 0$.

The goal is to develop effective methods for computing minimizers of such functionals over classes of functions of practical interest. Now this task, in even *simple* cases, such as the linear regression model with a quadratic penalty, is not completely trivial. When Q(x|z) is not quadratic these numerical difficulties can be greatly magnified. Our discussion will restrict itself to the cases where the penalty term, J, is quadratic in x. Extensions to uniformly positive definite penalties should be fairly obvious.

To begin with, many numerical minimization algorithms proceed by successively obtaining quadratic approximations to the objective function; the minima of these quadratics are used to define a sequence of iterates which, one hopes, will converge to a minimizer of the objective function. Due care has to be taken to insure that the successive quadratic approximants are computationally tractable, i.e. their minimizers are easily obtainable. This consideration is especially important for infinite dimensional problems. We restrict ourselves to fixed-step size Newton-like schemes. Iterates, $\{x^k\}$, are defined by:

$$x^{k+1} = x^k - [A(x^k)]^{-1} \Delta I_{\lambda}(x^k)$$
 $k=0,1,...$ (*)

where $A(x^k) \in L(H,H) \forall k$. Obvious choices for A come to mind. The standard Newton choice is:

$$A(x^{k}) = H_{I_{\lambda}}(x^{k})$$
$$= H_{Q}(x^{k}) + \lambda H_{J}(x^{k})$$

where H denotes the Hessian. With Fisher's Scoring Technique $A(x^k)$ becomes:

$$A(x^{k}) = E\{H_{I_{\lambda}}(x^{k})\}$$
$$= E\{H_{Q}(x^{k})\} + \lambda H_{J}(x^{k})$$

Here expectation is taken assuming x^k is the true value of x .

Finally if Q(x|z) = (F(x|z), F(x|z)) where (\cdot, \cdot) is some inner product on the range space of F, $A(x^k)$ might be chosen according to a Gauss-Newton formulation i.e.:

$$A(x^{k}) = (\nabla F(x^{k} | z), \nabla F(x^{k} | z)) + H_{J}(x^{k}).$$

Our reason for favouring schemes of the above type is that these methods can often be easily meshed with existing finite-dimensional software. The final section of this chapter outlines the implementation of such algorithms in Generalized Linear Interactive Models and Non-linear Regression Model contexts.

As with any iterative algorithm one is naturally interested in knowing convergence properties. Does the process converge to a minimizer of I_{λ} ? If so, what is the rate of convergence? Our goal in this chapter is to show that some basic results concerning the convergence of Newton-like schemes (*), which have already been established in finite-dimensional situations, can be generalized to handle the broad class of minimization problems associated with penalized likelihood methods. Results about convergence properties are presented in the next section. Some specifics on implementation of the techniques for Generalized Linear Models and Non-Linear Regression Models are outlined in section 3 of this chapter. In Chapter 5, a non-linear integral equation arising in Satellite Meteorology is numerically inverted using the algorithm discussed here - the method appears to work quite well in practice.

The major reference for this Chapter is Ortega and

Rheinbold (1970).

3.2 CONVERGENCE PROPERTIES OF THE ALGORITHMS

Asymptotic Root Convergence Factors

Before getting into the convergence analysis, we first need to introduce a way of measuring the asymptotic rate of convergence of a convergent sequence.

Definition 3.2.1.

Suppose $\{x_k\}$ is a sequence in H converging to a point x^* . The asymptotic root convergence factors, $R_p\{x^k\}$, of the sequence are given by:

$$\mathbb{R}_{p} \{ \mathbf{x}^{k} \} = \begin{cases} \lim_{k \to \infty} \sup_{k \to \infty} \|\mathbf{x}^{k} - \mathbf{x}^{\star}\|^{1/k} & \text{if } p=1 \\ \lim_{k \to \infty} \sup_{k \to \infty} \|\mathbf{x}^{k} - \mathbf{x}^{\star}\|^{1/p^{k}} & \text{if } p>1 \end{cases}$$

Notice that whenever $R_p\{x^k\} < 1$, then, for any $\varepsilon > 0$ with $R_p\{x^k\} + \varepsilon < 1$, there is a $k_0 \ge 0$, such that $\forall \ k \ge k_0$ either

$$||x^{k} - x^{*}|| \le (R_{p}\{x^{k}\} + \varepsilon)^{p^{k}}$$
 if $p > 1$

or

$$\|x^{k} - x^{*}\| \le (R_{1}\{x^{k}\} + \varepsilon)^{k}$$
 if p=1.

In the latter case, the convergence of the sequence is at least as rapid as a geometric progression with ratio

 $R_1 + \varepsilon < 1$. If $0 < R_1\{x^k\} < 1$, the convergence is called R-linear, while for $R_1\{x^k\} = 1$ or $R_1\{x^k\} = 0$ convergence is R-sublinear and R-superlinear, respectively. Similarly, if

$$0 < R_{2}\{x^{k}\} < 1$$

we say that the convergence is R-quadratic.

More detailed information about this measure of asymptotic rate of convergence can be found in chapter 9 or Ortega and Rheinbold (1970).

Convergence Analysis

We now turn to the analysis of the algorithms. Throughout this discussion the penalty term $J(x) = ||Px||^2$ where P is, as usual, a projection operator on H with finite-dimensional null space. The objective functionals therefore have the form:

$$I_{\lambda}(x) = Q(x) + \lambda \|Px\|^{2} \quad \lambda > 0 \text{ and } x \in H \quad (3.2.1)$$

Q will be assumed to be bounded below. We wish to find a minimizer of I_{λ} in a closed convex subset C of H. Let x^k be the current estimate and x^{k+1} the update. We consider iterative schemes of the form:

$$x^{k+1} = x^k - [A(x^k)]^{-1} \nabla I_{\lambda}(x^k) \qquad k=0,1,...$$
 (3.2.2)

where $A(x^k) \in L(H,H) \forall k$.

In this section we prove the following result about the convergence of iterations of the form (3.2.2.).

Theorem 3.2.2.

Suppose Q in equation (3.2.1) is such that Q" is continuous and Q' is weakly continuous on the interior of C. Let $x_{\lambda}^{0} \in int C$ be such that the level set

$$\mathbf{L}^{0} = \{\mathbf{x} \mid \mathbf{I}_{\lambda}(\mathbf{x}) \leq \mathbf{I}_{\lambda}(\mathbf{x}_{\lambda}^{0}) \}$$

is weakly compact and I_{λ} has only finitely many critical points in L^0 . Assume that $A:L^0 \rightarrow L(H,H)$ is a continuous mapping such that Ξ an invertible map $B(x) \in L(H,H)$ with $A(x) = B^*(x)B(x)$ (* denotes the adjoint) $\forall x \in L^0$. Suppose also $\Xi \mu_0$, μ_1 and γ_1 all positive with $0 < \frac{1}{2}\gamma_1 < .$ μ_0 satisfying

$$\mu_0 \|h\|^2 \le \langle h, A(x)h \rangle \le \mu_1 \|h\|^2$$

and

$$I_{1}''(\mathbf{x}) hh \leq \gamma_{1} ||h||^{2}$$

 $\forall x \in L^0$ and $h \in H$. Consider the process (3.2.2) i.e.

$$x^{k+1} = x^k - [A(x^k)]^{-1} \nabla I_{\lambda}(x^k) \qquad k=0,1,...$$

Then $\{x^k\} \in L^0$, $\lim_{k \to \infty} x^k = x^*$, where $\forall I_{\lambda}(x^*) = 0$, and if

 $H_{I_{\lambda}}(x^*)$ is nonsingular, then the convergence is at least R-linear.

Comment.

A finite-dimensional version of this Theorem appears in Chapter 14 of Ortega and Rheinbold (1970). The proof we give here parallels their argument. Care, however, is needed in dealing with the *weak* compactness of L^0 . The convergence properties of the Newton method (A(x) = $H_Q(x) + \lambda H_J(x)$) in abstract Banach spaces have, of course, been known for some time. The famous Newton-Kantorovic Theorem (see Chapter 4 of Rall (1969)) gives conditions under which the Newton method converges R-quadratically. In view of this one might expect that the rate of convergence for Fisher's scoring technique is better than R-linear.

We begin the proof of Theorem 3.2.2 by establishing four lemmas which will be used in the proof.

Lemma 3.2.3.

For $\lambda > 0$ let $I_{\lambda} = Q(x) + \lambda \|Px\|^2$, $x \in D \subset H$ where $Q: D \to R^1$ has a weakly continuous first Frechet derivative on $D_0 - D_0$, a weakly compact subset of D. Suppose that $\{x^k\} \in D_0$ is any sequence which satisfies $\lim_{k \to \infty} \|\nabla I_{\lambda}(x^k)\| = 0$. Then the set $\Omega = \{x \in D_0 | \nabla I_{\lambda}(x) = 0\}$ of critical points of I_{λ} in D_0 is non-empty and

$$\lim_{k \to \infty} \inf_{x \in \Omega} \|x^{k} - x\| = 0.$$

In particular if Ω consists of a single point x^* Then $\lim_{k \to \infty} x^k = x^*$ and $\nabla I_{\lambda}(x^*) = 0$.

Proof:

Since D_0 is weakly compact $\{x^k\}$ has a weakly convergent subsequence $\{x^{ki}\}$ with $x^{ki} \stackrel{W}{\rightarrow} x^*$ for some $x^* \in D_0$. Now $\nabla I_{\lambda}(x^{ki}) = \nabla Q(x^{ki}) + 2\lambda P x^{ki}$ and so by the weak continuity of Q' on $D_0 \quad \nabla I_{\lambda}(x^{ki}) \stackrel{W}{\rightarrow} \nabla I_{\lambda}(x^*)$. Therefore, $|\langle \nabla I_{\lambda}(x^*), \nabla I_{\lambda}(x^{ki}) \rangle| \rightarrow ||\nabla I_{\lambda}(x^*)||^2$. However $|\langle \nabla I_{\lambda}(x^*), \nabla I_{\lambda}(x^{ki}) \rangle| \leq ||\nabla I_{\lambda}(x^*)|| \quad ||\nabla I_{\lambda}(x^{ki})|| \quad (Cauchy-Schwartz Inequality)$ and since $||\nabla I_{\lambda}(x^{ki})|| \rightarrow 0$ (by definition of $\{x^{ki}\}$) we have $\nabla I_{\lambda}(x^*) = 0$. That is to say Ω is non-empty. But this also means that

$$\lim_{k_{i} \to \infty} \| \nabla I_{\lambda}(x^{k_{i}}) - \nabla I_{\lambda}(x^{*}) \| = 0.$$

Now

$$\nabla I_{\lambda}(\mathbf{x}^{k_{\mathbf{j}}}) - \nabla I_{\lambda}(\mathbf{x}^{*}) = \nabla Q(\mathbf{x}^{k_{\mathbf{j}}}) - \nabla Q(\mathbf{x}^{*}) + 2\lambda (P\mathbf{x}^{k_{\mathbf{j}}} - P\mathbf{x}^{*})$$

which implies

$$\begin{aligned} & 2\lambda \| \mathbf{P} \mathbf{x}^{\mathbf{k}\mathbf{i}} - \mathbf{P} \mathbf{x}^{*} \| = \| \nabla \mathbf{I}_{\lambda} (\mathbf{x}^{\mathbf{k}\mathbf{i}}) - \nabla \mathbf{I}_{\lambda} (\mathbf{x}^{*}) - \nabla \mathbf{Q} (\mathbf{x}^{\mathbf{k}\mathbf{i}}) + \nabla \mathbf{Q} (\mathbf{x}^{*}) \| \\ & \leq \| \nabla \mathbf{I}_{\lambda} (\mathbf{x}^{\mathbf{k}\mathbf{i}}) - \nabla \mathbf{I}_{\lambda} (\mathbf{x}^{*}) \| + \| \nabla \mathbf{Q} (\mathbf{x}^{\mathbf{k}\mathbf{i}}) - \nabla \mathbf{Q} (\mathbf{x}^{*}) \| \end{aligned}$$

Hence, recalling that Q' is weakly continuous, we have

$$Px^{K_{i}} \rightarrow Px^{*}$$
 as $i \rightarrow \infty$.

But since the null space of P is finite-dimensional and $x^{k_i} \xrightarrow{W} x^*$ we have that

$$(I - P)x^{k_{i}} \rightarrow (I - P)x^{*}$$
 as $i \rightarrow \infty$.

Hence $x^{k_{i}} \rightarrow x^{*}$ as $i \rightarrow \infty$. To finish off the proof let $\delta_{k} = \inf_{\substack{x \in \Omega \\ x \in \Omega}} \|x^{k} - x\|$ and suppose \exists a subsequence $\{\delta_{k_{j}}\}$ such that $\delta_{k_{j}} \ge \delta > 0 \quad \forall j$. Then, by an argument similar to that used above, the corresponding sequence of x's, $\{x^{k_{j}}\}$ say, must have a convergent subsequence whose limit lies in Ω . But this contradicts the definition of $\delta_{k_{j}}$. Therefore $\lim_{k \rightarrow \infty} \delta_{k} = 0$.

The proof of the next two lemmas is practically the same as the proofs of (14.1.15) and (14.1.16) in Ortega and Rheinbold (1970). They are included for completeness.

Lemma 3.2.4.

Let I_{λ} , D_0 and D be as in Lemma 3.2.3. Suppose that the set Ω of critical points of I_{λ} is D_0 is finite. Let $\{x^k\} \in D_0$ be any sequence for which $\lim_{k \to \infty} ||x^k - x^{k+1}|| = 0$ and $\lim_{k \to \infty} ||\nabla I_{\lambda}(x^k)|| = 0$. Then $\lim_{k \to \infty} x^k = x^*$ and $\nabla I_{\lambda}(x^*) = 0$.

Proof:

Let Γ be the set of limit points of $\{x^k\}$. As in

Lemma 3.2.3, any limit point is also a critical point of I_{λ} so that $\Gamma \subset \Omega$; that is, Γ is finite. Suppose that $\Gamma = \{z^1, z^2, \dots, z^m\}$ with m > 1; then

$$\delta = \min \{ \|z^{i} - z^{j}\| : i \neq j, i, j = 1, ..., m \} > 0$$

and we can choose $k_0 \ge 0$ such that $x^k \in \bigcup_{i=1}^m S(z^i, \delta/4)$ and $\|x^k - x^{k+1}\| \le \delta/4$ for all $k \ge k_0$. $x^k \in S(z^1, \delta/4)$ implies that

$$\|z^{i} - x^{k+1}\| \ge \|z^{i} - z^{1}\| - (\|z^{1} - x^{k}\| + \|x^{k} - x^{k+1}\|)$$

 $\geq \delta - 2\delta/4 = \delta/2, \quad i \geq 2$

and hence, necessarily, that $x^{k+1} \in S(z^1, \delta/4)$. By induction, therefore, $x^k \in S(z^1, \delta/4)$ for all $k \ge k_0$, which contradicts the fact that z^2, \ldots, z^m are limit points of $\{x^k\}$; therefore, m = 1. \Box

Lemma 3.2.5.

Let $g: D \subset \mathbb{R}^n \to \mathbb{R}^1$ be Frechet differentiable on an open set $D_0 \subset D$ and suppose that $\{x^k\} \subset D_0$ converges to $x^* \in D_0$. Assume that $\forall g(x^*) = 0$, and g has a second Frechet derivative at x^* and the Hessian $H_g(x^*)$ is invertible, and that there is an $\eta > 0$ and a k_0 for which

$$g(x^{k}) - g(x^{k+1}) \ge \eta \|\nabla g(x^{k})\|^{2}, \quad \forall k \ge k_{0}.$$
 (3.2.3)

Then the root convergence factor $R_{l}\{x^{k}\} < 1$. That is, the sequence converges, at least, R-linearly.

Proof:

For $A \in L(H,H)$

$$\|A\| \equiv \sup_{\mathbf{y} \in \mathbf{H}} \frac{|\langle \mathbf{y}, A\mathbf{y} \rangle|}{\|\mathbf{y}\|^2} .$$

Set $\alpha = \|H_{g}(x^{*})^{-1}\|^{-1}$. Then for given $\varepsilon \in (0, \alpha)$ we may choose $\delta > 0$ so that $S \equiv S(x^{*}, \delta) \in D_{0}$ and

$$\|\nabla g(\mathbf{x}) - H_{\alpha}(\mathbf{x}^{*})(\mathbf{x} - \mathbf{x}^{*})\| \leq \varepsilon \|\mathbf{x} - \mathbf{x}^{*}\|, \quad \forall \mathbf{x} \in S.$$

Hence

$$\begin{aligned} \|\nabla g(x)\| \ge \|H_{g}(x^{*})(x - x^{*})\| - \|\nabla g(x) - H_{g}(x^{*})(x - x^{*})\| \\ \ge (\alpha - \varepsilon)\|x - x^{*}\|, \quad \forall \ x \in S . \end{aligned}$$

Consequently in view of equation (3.2.3), there is a k_0 so that with $\gamma_0 = [\eta (\alpha - \varepsilon)^2]^{-1}$

$$\|\mathbf{x}^{k} - \mathbf{x}^{*}\| \leq (\alpha - \varepsilon)^{-2} \|\nabla g(\mathbf{x}^{k})\|^{2}$$
$$\leq \gamma_{0} [g(\mathbf{x}^{k}) - g(\mathbf{x}^{k+1})], \quad \forall \ k \geq k_{0}. (3.2.4)$$

Now note that the mean-value Theorem 1.A.8 ensures that

$$g(x^{k}) - g(x^{*}) \leq \gamma_{1} \|x^{k} - x^{*}\|^{2}, \forall k \geq k_{2},$$

where $\gamma_1 = \frac{1}{2} \|H_g(x^*)\| + \varepsilon$, and $k_2 \ge k_0$ is sufficiently large. Hence, we obtain, using 3.2.3 and 3.2.4,

$$0 \leq g(x^{k+1}) - g(x^{*}) \leq g(x^{k}) - g(x^{*}) - \eta \|\nabla g(x^{k})\|^{2}$$

$$\leq g(x^{k}) - g(x^{*}) - (1/\gamma_{0}) \|x^{k} - x^{*}\|^{2}$$

$$\leq \mu[g(x^{k}) - g(x^{*})], \quad \forall \ k \geq k_{2}, \qquad (3.2.5)$$

where $\mu = 1 - (\gamma_0 \gamma_1)^{-1}$. Clearly (by a suitable choice of η) $0 \le \mu < 1$. Now letting $e_k = [g(x^k) - g(x^*)]^{\frac{1}{2}}$ for $k = 1, 2, \ldots$, we have for $k \ge k_2$

so by induction

$$e_{k} \leq \mu^{(k-k_{2})/2} e_{k_{2}} \leq \mu^{(k-k_{2})/2} [1+e_{k_{2}}].$$

Therefore

$$(e_k)^{1/k} \leq \mu^{(k-k_2)/2k} [1+e_{k_2}]^{1/k}$$

which implies

$$\lim_{k} \sup_{k} (e_{k})^{1/k} \leq \lim_{k} \sup_{k} \mu^{(k-k_{2})/2k} [1+e_{k_{2}}]^{1/k}$$
$$\leq \mu^{\frac{1}{2}} < 1.$$

That is $R_1\{[g(x^k) - g(x^*)]^{\frac{1}{2}}\} < 1$. Therefore, by (3.2.4) and remembering that $g(x^k) - g(x^{k+1}) \le g(x^k) - g(x^*)$,

$$R_{1}\{x^{k}\} \leq R_{1}\{[g(x^{k}) - g(x^{*})]^{\frac{1}{2}}\} < 1.$$

Lemma 3.2.6.

Let $g: D \subset H \rightarrow R$ be continuous on an open set D and Frechet differentiable on

$$L^{0} = L^{0}[g(x^{0})] = \{x \in D | g(x) \le g(x^{0})\}$$

for some x^0 in D. Suppose L^0 is weakly compact. Then for any $x \in L^0$ and $p \in H$ with $\langle \nabla g(x), p \rangle > 0$ there is an $\alpha^* > 0$ such that

$$g(x) = g(x - \alpha^* p)$$
 and $[x, x - \alpha^* p] \subset L^0$.

In particular if $\eta > 0$ is any number with the property that

$$q(x-\alpha p) < q(x) \quad \forall x-\alpha p \in (x, x-\eta p] \cap L^{\circ}$$

then $[x, x - \eta p] \subset L^0$.

Proof:

Set $\alpha^* = \sup J$ where

 $J = \{\alpha > 0 \mid [x, x - \alpha p] \subset D \text{ and } g(x - \beta p) < g(x), \forall \beta \in (0, \alpha] \}.$

The function $\zeta(\alpha) = g(x - \alpha p) - g(x)$, $\alpha \in \mathbb{R}^+$ is continuous, $\zeta(0) = 0$ and $\zeta'(0) = -g'(x)p = -\langle \nabla g(x), p \rangle < 0$. Hence $\exists \alpha > 0$ s.t. $\forall \beta \in (0, \alpha], \zeta(\beta) < 0$. It follows from this and the openness of D that J is non-empty and so α^* is well-defined. Consider the half line $\ell = \{x - \alpha p, \alpha > 0\}$. Since L^0 is weakly compact, $L^0 \cap \ell$ is compact. Therefore the set $\{\alpha > 0: x - \alpha p \in L^0\}$ must also be compact. Which of course means that $\alpha^* < \infty$ and $[x, x - \alpha^* p] \in L^0$.

Now suppose that $g(x - \alpha^* p) < g(x)$. Here since D is open and g is continuous $\exists \delta > 0$ such that $x - \alpha p \in D$ and $g(x - \alpha p) < g(x) \forall \alpha \in [\alpha^*, \alpha^* + \delta]$. But this contradicts the definition of α^* and so $g(x - \alpha^* p) = g(x)$.

The last statement follows immediately since if $\eta \ge \alpha^*$ then $g(x - \alpha^*p) < g(x)$ which contradicts the definition of α^* .

We are now ready to prove Theorem 3.2.2. Recall the statement.

Theorem 3.2.2.

Suppose Q in equation (3.2.1) is such that Q" is continuous and Q' is weakly continuous on the interior of C. Let $x_{\lambda}^{0} \in \text{int C}$ be such that the level set

$$\mathbf{L}^{0} = \{ \mathbf{x} \in \mathbf{C} \, \big| \, \mathbf{I}_{\lambda}(\mathbf{x}) \leq \mathbf{I}_{\lambda}(\mathbf{x}_{\lambda}^{0}) \}$$

is weakly compact and I_{λ} has only finitely many critical points in L^0 . Assume that A: $L^0 \rightarrow L(H,H)$ is a continuous

mapping such that I an invertible map $B(x) \in L(H,H)$ with $A(x) = B^*(x)B(x)$ (* denotes the adjoint) $\forall x \in L^0$. Suppose also μ_0 , μ_1 and γ_1 all positive with $0 < \frac{1}{2}\gamma_1 < \mu_0$ satisfying

$$\mu_0 \|h\|^2 \le \langle h, A(x)h \rangle \le \mu_1 \|h\|^2$$

and

$$I_{\lambda}^{"}(x) hh \leq \gamma_{1} \|h\|^{2}$$

 $\forall x \in L^0$ and $h \in H$. Consider the process (3.2.2) i.e.

$$x^{k+1} = x^k - [A(x^k)]^{-1} \nabla I_{\lambda}(x^k) \qquad k = 0, 1, \dots$$

Then $\{x^k\} \in L^0$, $\lim_{k \to \infty} x^k = x^*$, where $\forall I_{\lambda}(x^*) = 0$, and if $H_{I_{\lambda}}(x^*)$ is nonsingular, then the convergence is at least R-linear.

Proof:

For ease of notation let $g \equiv I_{\lambda}$ i.e.

$$q(x) = Q(x) = \lambda ||Px||^2.$$

For any $x \in L^0$ and $h \in H$

$$\mu_0 \|h\|^2 \le \langle h, A(x) h \rangle \le \mu_1 \|h\|^2 . \qquad (3.2.6)$$

If $x^k \in L^0$, $B(x^k)$ is invertible, so

$$p^{k} = A(x^{k})^{-1} \nabla g(x^{k}) = B^{*}(x^{k})^{-1} B(x^{k})^{-1} \nabla g(x^{k})$$

is well-defined. Letting $h = B(x^k)^{-1} \nabla g(x^k)$ we have from (3.2.6) that

$$< B(x^{k})^{-1} \nabla g(x^{k}), A(x^{k}) B(x^{k}) \nabla g(x^{k}) > \leq \mu_{1} < B(x^{k})^{-1} \nabla g(x^{k}), B(x^{k})^{-1} \nabla g(x^{k}) >$$

which by the Cauchy-Schwartz inequality implies that $\frac{1}{\mu_1} \| \nabla g(x^k) \| \le \| p^k \|$. Similarly letting $h = A(x^k)^{-1} \nabla g(x^k)$ in (3.2.1) we obtain that $\| p^k \| \le \frac{1}{\mu_0} \| \nabla g(x^k) \|$. Hence

$$(\mu_1)^{-1} \| \nabla g(\mathbf{x}^k) \| \le \| p^k \| \le (\mu_0)^{-1} \| \nabla g(\mathbf{x}^k) \|$$
 (3.2.7)

It is also clear that $\langle p^k, \nabla g(x^k) \rangle > 0$ unless $\|A(x^k)^{-1}\nabla g(x^k)\| = 0$ i.e. $\nabla g(x^k) = 0$. In other words, $\langle p^k, \nabla g(x^k) \rangle > 0$ unless the process terminates at x^k . Now let $\alpha \in (0,1]$ be such that $[x^k, x^k - \alpha p^k] \subset L^0$. Then, by the mean-value Theorem 1.A.7, there is an $\hat{\alpha} \in (0,\alpha]$ so that

$$g(x^{k}) - g(x^{k} - \alpha p^{k}) = \alpha \langle \forall g(x^{k}), p^{k} \rangle - \frac{1}{2}\alpha^{2} \langle p^{k}, H_{g}(x^{k} - \hat{\alpha} p^{k}) p^{k} \rangle$$

But since $\langle \nabla g(x^k), p^k \rangle = \langle p^k, A(x^k)p^k \rangle$ and by hypothesis, $\langle p^k, H_g(x^k - \hat{\alpha}p^k)p^k \rangle \leq \frac{1}{2} \gamma_1 \|p^k\|^2$ since $x^k - \hat{\alpha}p^k \in L^0$. $\therefore g(x^k) - g(x^k - \alpha p^k) = \alpha [\mu_0 - \frac{1}{2} \gamma_1] \|p^k\|^2 > 0$. (3.2.8) It follows from Lemma 3.2.2 that $x^{k+1} = x^k - p^k \in L^0$ and,

by induction, that $\{x^k\} \in L^0$.

Combining (3.2.7) and (3.2.8) we get

$$g(x^{k}) - g(x^{k+1}) \ge (\mu_{1})^{-2}(\mu_{0} - \frac{1}{2}\gamma_{1}) \|\nabla g(x^{k})\|^{2}$$
 (3.2.9)

So, since $\{g(x^k)\}$ is a monotonic decreasing sequence which is bounded below, the left hand side converges to zero. Hence $\lim_{k \to \infty} \nabla g(x^k) = 0$. Consequently, since $p^k = x^k - x^{k+1}$, (3.2.7) shows that $\lim_{k \to \infty} (x^{k+1} - x^k) = 0$, and the result follows by Lemmas 3.2.4 and 3.2.5.

3.3 SPECIFIC IMPLEMENTATIONS OF THE ALGORITHM

We now give some details on the implementation of iterative minimization methods of the form

$$\mathbf{x}^{k+1} = \mathbf{x}^k - [\mathbf{A}(\mathbf{x}^k)]^{-1} \nabla \mathbf{I}_{\mathbf{x}}(\mathbf{x}^k)$$

in Generalized Linear Interative Models and Non-linear Regression models. Let x^0 be the current iterate. The following describes how the next iterate, x^1 , is obtained. We begin by describing a Newton-Raphson method for use with Generalized Linear Interactive Models.

GENERALIZED LINEAR INTERACTIVE MODELS

Recall from Chapter 1 that I_{λ} has the form:

$$I_{\lambda}(\mathbf{x}) = \sum_{i=1}^{n} [b(\theta_{i}) - y_{i}\theta_{i}]/a_{i}(\phi) + \lambda J(\mathbf{x}) = Q(\mathbf{x}|\mathbf{y}) + \lambda J(\mathbf{x})$$

We will show that minimizing the Newton-Raphson approximation (with Fisher's Scoring technique)

$$I_{\lambda}(x) \approx Q(x^{0}) + Q'(x^{0})(x - x^{0}) + \frac{1}{2}E(Q''(x^{0}))(x - x^{0})(x - x^{0}) + \lambda J(x)$$
(*)

is equivalent to minimizing I^0_λ given by:

$$I_{\lambda}^{0}(\mathbf{x}) = \sum_{i=1}^{n} \frac{1}{u_{i}} [z_{i} - L_{i}x]^{2} + \lambda J(\mathbf{x})$$

where

$$z_{i} = L_{i}x^{0} + (y_{i} - \mu_{i}) \frac{\partial \eta_{i}}{\partial \mu i}; u_{i} = 2 \text{ var } (y_{i}) (\frac{\partial \eta_{i}}{\partial \mu_{i}})^{2}$$

One should note from this that the determination of the next iterate, x^1 , reduces to a generalized spline smoothing problem (see Wendelberger 1982). The derivation essentially follows Nelder and Wedderburn (1972).

Let the contribution of y_i to I_λ be denoted l_i

$$\ell_{i} = [b(\theta_{i}) - y_{i}\theta_{i}]/a_{i}(\phi) .$$

Recall the following relations:

$$E(y_i) = b'(\theta_i) = \mu_i, var(y_i) = b''(\theta_i)a_i(\theta) = \frac{\partial \mu_i}{\partial \theta_i}a_i(\theta)$$

and

$$\eta_i = L_i x = g(\mu_i) .$$

(See Chapter 1 for more details.) By the chain rule, the Frechet derivative of ℓ_i w.r.t. x, ℓ'_i is given by

$$\begin{aligned} \boldsymbol{x}_{i}^{\prime} &= \frac{\partial \boldsymbol{x}_{i}}{\partial \boldsymbol{\mu}_{i}} \boldsymbol{\mu}_{i}^{\prime} \\ &= \frac{\left[\boldsymbol{\mu}_{i} - \boldsymbol{y}_{i}\right]}{\operatorname{var}(\boldsymbol{y}_{i})} \quad \frac{\partial \boldsymbol{\mu}_{i}}{\partial \boldsymbol{\eta}_{i}} \boldsymbol{L}_{i} \end{aligned}$$

Similarly the second Frechet derivative is

$$\ell''_{i} = \frac{\partial \ell'_{i}}{\partial \mu_{i}} \frac{\partial \mu_{i}}{\partial \eta_{i}} L_{i}$$

$$= \left[\frac{\partial^{2} \ell_{i}}{\partial \mu_{i}^{2}} \frac{\partial \mu_{i}}{\partial \eta_{i}} L_{i} + \frac{\partial \ell_{i}}{\partial \mu_{i}} \frac{\partial^{2} \mu_{i}}{\partial \eta_{i} \partial \mu_{i}} L_{i}\right] \frac{\partial \mu_{i}}{\partial \eta_{i}} L_{i}$$

Using the fact that $E\left(\frac{\partial k_{i}}{\partial \mu_{i}}\right) = 0$ we obtain

$$E[\ell''_{i}] = E\left[\frac{\partial^{2}\ell_{i}}{\partial\mu_{i}^{2}}\right]\left(\frac{\partial\mu_{i}}{\partial\eta_{i}}\right)^{2} L_{i}L_{i}.$$

and $E\left[\frac{\partial^2 \ell_i}{\partial \mu_i^2}\right] = E\left[\frac{\partial}{\partial \mu_i}\left[\frac{\mu_i - y_i}{\text{var }(y_i)}\right]\right]$ which in turn is equal

to $1/var(y_i)$. By the preceding remarks, we have that

$$Q'(x^{0})(x - x^{0}) = \sum_{i=1}^{n} \ell_{i}'(x - x^{0})$$
$$= \sum_{i=1}^{n} \frac{\mu_{i} - y_{i}}{\operatorname{var}(y_{i})} \frac{\partial \mu_{i}}{\partial \eta_{i}} L_{i}(x - x^{0})$$

53

and also that

$$E[Q''(x^{0})](x - x^{0})(x - x^{0}) = \sum_{i=1}^{n} E[l_{i}''(x^{0})](x - x^{0})(x - x^{0})$$
$$= \sum_{i=1}^{n} \left(\frac{\partial \mu_{i}}{\partial \eta_{i}}\right)^{2} \frac{[L_{i}(x - x^{0})]^{2}}{var(y_{i})}$$

Combining terms it follows that minimizing (*) w.r.t. x is equivalent to minimizing I^0_λ given by

$$I_{\lambda}^{0}(\mathbf{x}) = \sum_{i=1}^{n} \frac{1}{u_{i}} [z_{i} - L_{i}x]^{2} + \lambda J(\mathbf{x})$$

where z_i and u_i are as given.

NON-LINEAR REGRESSION MODELS

In these models we have that

$$I_{\lambda}(\mathbf{x}) = \sum_{i=1}^{n} [z_i - N_i(\mathbf{x})]^2 + \lambda J(\mathbf{x}).$$

 N_i 's are non-linear functionals of x contaminated by noise. In the Gauss-Newton algorithm N_i 's are linearized about the current iterate as

$$N_{i}(x) \approx N_{i}(x^{0}) + N_{i}'(x^{0})(x - x^{0})$$

and the next iterate, x^1 , is the minimizer of the functional I^0_λ given by:

$$I_{\lambda}^{0}(x) = \sum_{i=1}^{n} [z_{i} - N_{i}(x^{0}) - N_{i}'(x^{0})(x - x^{0})]^{2} + \lambda J(x)$$

Again one should note from this that the determination of the next iterate, x^1 , reduces to a generalized spline smoothing problem.

CHAPTER 4

ESTIMATION OF THE SMOOTHING PARAMETER

4.1 INTRODUCTION

Recall that the penalized likelihood estimator, $\mathbf{x}_{\lambda'}$ is the minimizer of the functional

 $I_{\lambda}(x) = Q(x | z) + \lambda J(x) \quad \lambda > 0 .$

The *smoothing* parameter λ controls the relative weighting given to the penalty term, large values of λ giving the penalty more weight.

There is a growing literature on methods for estimating smoothing parameters (see Atilgan 1983, Anderssen and Bloomfield 1974, Leonard 1982, Silverman 1978, and Stone 1974), one of the most popular of these methods is crossvalidation. The cross-validatory procedure tries to choose a value of the smoothing parameter which is optimal for some specified predictive criterion (see Stone 1974). Roughly, the method works as follows: Data are divided into subsamples; for a given value of the smoothing parameter the estimator is computed on each of the subsamples and, using the predictive criterion, its performance is evaluated on the remaining subsamples. The "cross-validatory" assessment of λ is the combined performance of the estimator over all the subsamples. Using this technique the value of the smoothing parameter with the best cross-validatory assessment can be isolated.

Since Stone's 1974 paper, cross-validation has been widely applied in a number of diverse statistical contexts. Wahba and Wold (1975) were the first to use the method in a penalized likelihood setting. However, the method they proposed was computationally intensive and in 1979 Craven and Wahba came up with a refined procedure, Generalized Cross Validation, which has decided computational advantages over the Wahba and Wold proposal.

In this chapter we extend the Generalized Cross Validation (GCV) estimator to general penalized likelihood functionals. The estimator is described in the next section and some favourable heuristics are given in section 3. Although the results are of a preliminary nature, they suggest that the technique may have promise.

4.2 THE GENERALIZED CROSS VALIDATION PROCEDURE

Generalized Cross-Validation (GCV) was developed by Craven and Wahba (1979) to estimate a smoothing parameter similar to λ in a linear regression model context. In their model, the observations z_i are noise contaminated linear functionals of an unknown function, x, of interest.

$$z_{i} = L_{i}x + noise, \quad i = 1, 2, ..., m$$

where the L_i 's are linear functionals. The estimator, x_{λ} , of x is the minimizer of the functional

$$\frac{1}{n} \sum_{i=1}^{n} [z_i - L_i x]^2 + \lambda \int [x^m(t)]^2 dt \quad \lambda > 0 .$$

It turns out that the residuals, $z_i - \hat{z}_i$, can be written as a linear transform of the data.

$$z - \hat{z} = [I - A(\lambda)]z$$

where $A(\lambda)$ is an n n matrix defined in terms of the linear functionals. Formally $A(\lambda)_{ij} = \frac{\partial L_i}{\partial z_j}(x_{\lambda})$. The GCV estimate of λ is the minimizer of the function

$$V(\lambda) = \frac{\frac{1}{n} \sum_{i=1}^{n} [z_i - \hat{z}_i]^2}{(\frac{1}{n} \operatorname{tr} [I - A(\lambda)])^2}$$

$$V(\lambda) = \frac{\frac{1}{n} RSS(\lambda)}{[1-\mu_1(\lambda)]^2}$$

 $V(\lambda)$ is an assessment of λ combining goodness of fit and model complexity. The goodness of fit is measured by the residual sum of squares, $\frac{1}{n} RSS(\lambda)$, while the complexity of x_{λ} is measured by the "effective" number of parameters describing x_{λ} relative to the sample size, i.e. $\mu_{1}(\lambda)$. (See Wahba (1979) for an equivalent degreesof-freedom interpretation of tr[I-A(λ)].) As $\lambda \neq \infty$, x_{λ} tends to an m-1'th degree polynomial regression making $\frac{1}{n}$ RSS(λ) take on its maximum and $\mu_{1}(\lambda)$ be a minimum. On the other hand, as $\lambda \neq 0$, x_{λ} tends to a data interpolant minimizing $\frac{1}{n}$ RSS(λ) and maximizing $\mu_{1}(\lambda)$. It has been shown that the GCV estimate of λ , for large sample sizes, comes close to minimizing the integrated squared error, R(λ), between x and x_{λ} :

$$R(\lambda) = \frac{1}{n} \sum_{i=1}^{n} [L_{i}x - L_{i}x_{\lambda}]^{2}$$

For further discussion see articles by Craven and Wahba (1979) and Speckman (1982). Wahba (1982) proposed a GCV based estimator of λ in a constrained version of the linear regression model. In order to extend the GCV procedure to general penalized likelihoods, we need to introduce some additional notation.

Some Notation and Mild Assumptions

Let x_0 be the true value of the function, i.e. the value corresponding to the observed data z. Define the following information operators on H by:

i.1
$$I_{sp}(x_0) = EH_Q[x_0] + \lambda H_J[x_0]$$

i.2 $I_s(x_0) = EH_Q[x_0]$

 $\rm H_Q$ is the Hessian of Q, $\rm H_J$ the Hessian of J and expectations are taken with respect to the probability distribution of z.

We assume that $I_s(x_0)$ and $I_{sp}(x_0)$ (and their inverses) have cholesky-like decompositions (i.e. $I_s(x_0) = I_s^{\frac{1}{2}}(x_0)I_s^{\frac{1}{2}}(x_0)$) and that Q satisfies:

q.l	$E\left[\nabla Q\left(x_{0} \mid z\right)\right] \equiv 0$
q.2	$Var[\nabla Q(x_0 z)] = I_s(x_0)$
q.3	$\mathbb{E} \ \mathbf{I}_{s}^{-\frac{1}{2}}(\mathbf{x}_{0}) \nabla Q(\mathbf{x}_{0} \mathbf{z}) \ ^{2} = v_{0} < \infty$

where the inverses are generalized inverses.

Finally let $\mu_1(\lambda) = \nu_1(\lambda)/\nu_0$ and $\mu_2(\lambda)/\nu_0$ where ν_1 and ν_2 are given by:

$$\begin{aligned} v_{1}(\lambda) &= E \| \mathbf{I}_{sp}^{-\frac{1}{2}} \nabla Q(\mathbf{x}_{0} | z) \|^{2} \\ &= tr \{ \mathbf{I}_{s}(\mathbf{x}_{0}) \mathbf{I}_{sp}^{-1}(\mathbf{x}_{0}) \} \\ v_{2}(\lambda) &= E \| \mathbf{I}_{s}^{\frac{1}{2}}(\mathbf{x}_{0}) \mathbf{I}_{sp}^{-1}(\mathbf{x}_{0}) \nabla Q(\mathbf{x}_{0} | z) \|^{2} \\ &= tr \{ \mathbf{I}_{s}^{2}(\mathbf{x}_{0}) \mathbf{I}_{sp}^{-2}(\mathbf{x}_{0}) \} . \end{aligned}$$

(We tacitly assume the commutivity of the I-operators.)

The Proposed GCV Estimator

In the above notation, Craven and Wahba's GCV estimator is the minimizer of:

$$V(\lambda) = \frac{Q(x_{\lambda} | z)}{[1 - \mu_{1}(\lambda)]^{2}}$$

where x_{λ} is the minimizer of I_{λ} , and $Q(x_{\lambda}|z)$ is given by:

$$Q(x_{\lambda} | z) = \frac{1}{n} \sum_{i=1}^{n} [z_i - L_i x_{\lambda}]^2$$
.

This suggests that a plausible GCV estimator for general Penalized Likelihoods is

$$\nabla(\lambda) = \frac{Q(\mathbf{x}_{\lambda} | \mathbf{z})}{[1 - \mu_{1}(\lambda)]^{2}}$$

where μ_1 is now evaluated at x_{λ} ;

$$\mu_{l}(\lambda) = tr[I_{s}(x_{\lambda})I_{sp}^{-1}(x_{\lambda})]/\nu_{0}.$$

The matrices required to calculate μ_1 are usually at the final stage of the iterative procedure used to calculate x_{λ} (see Chapter 5 for an example).

4.3 SOME JUSTIFICATION FOR USING THE GCV ESTIMATOR

We now give some heuristic arguments in favour of the

GCV estimator of λ proposed above. The argument relies heavily on how well x_{λ} can be approximated by x_{λ}^{0} :

$$\mathbf{x}_{\lambda}^{0} = \mathbf{x}_{0} - \mathbf{I}_{sp}^{-1}(\mathbf{x}_{0}) \nabla \mathbf{I}_{\lambda}(\mathbf{x}_{0})$$

 x_{λ}^{0} also comes up in the asymptotic analysis of $x_{\lambda}^{}$. Loosely, one can show (under suitable regularity conditions) that if $\lambda \neq 0$ at an appropriate rate, the distance between x_{λ}^{0} and x_{0} is o_{p} of distance between x_{λ} and x_{λ}^{0} (see Cox and O'Sullivan 1983). By strengthening such results it may be possible to rigourize the heuristics which follow.

Suppose we say that:

$$EQ(x_{\lambda} | z) = EQ(x_{0} | z) + E < \nabla Q(x_{0} | z), x_{\lambda}^{0} >$$
$$+ \frac{1}{2}E < I_{s}(x_{0})(x_{\lambda} - x_{0}), (x_{\lambda} - x_{0}) >$$

and

$$\mu_{1}(\lambda) = tr[I_{s}(x_{0})I_{sp}^{-1}(x_{0})]/\nu_{0}$$

Then $2EV(\lambda)$ can be written as:

$$2\text{EV}(\lambda) = \frac{2\text{EQ}(\mathbf{x}_0 \mid \mathbf{z}) - 2\mu_1(\lambda)\nu_0 + \text{ER}(\lambda)\nu_0}{[1 - \mu_1(\lambda)]^2}$$

in which $ER(\lambda)$

$$ER(\lambda) = E < I_{s}(x_{0}) (x_{\lambda} - x_{0}), (x_{\lambda} - x_{0}) > /v_{0}$$
$$= b^{2}(x_{\lambda}) + v_{2}(x_{\lambda})$$
$$= higs + variability$$

where

$$b^{2}(x_{\lambda}) = \langle I_{s}(x_{0}) (Ex_{\lambda} - x_{0}) , (Ex_{\lambda} - x_{0}) \rangle / v_{0}$$

and

$$v_2(x_{\lambda}) = E \langle I_s(x_0)(x_{\lambda} - Ex_{\lambda}), (x_{\lambda} - Ex_{\lambda}) \rangle / v_0$$

Notice how the information operator enters naturally into ER(λ). For example, in the logistic regression model of section 3 of Chapter 2, R(λ) takes the form:

$$R(\lambda) = \frac{1}{n} \sum_{i=1}^{n} n_{i} p_{i} (1 - p_{i}) [L_{i} x_{\lambda} - L_{i} x_{0}]^{2}.$$

Assuming $v_2(x_{\lambda}) \approx v_2(x_{\lambda}^0) = \mu_2(\lambda)$ we get (using Craven and Wahba arguments) that:

$$\frac{\text{ER}(\lambda) + 2\text{EQ}(x_0 | z) - 2\text{EV}(\lambda)}{\text{ER}(\lambda)} = \frac{-\mu_1 (2 - \mu_1)}{(1 - \mu_1)^2} + \frac{1}{b^2 + \mu_2} - \frac{\mu_1^2}{(1 - \mu_1)^2}$$

from which it follows that

$$\frac{\mathrm{ER}(\lambda) + 2\mathrm{EQ}(\mathbf{x}_0 | \mathbf{z}) - 2\mathrm{EV}(\lambda)|}{\mathrm{ER}(\lambda)} < \mathrm{h}(\lambda)$$

where

$$h(\lambda) = \left[2\mu_{1}(\lambda) + \frac{\mu_{1}^{2}(\lambda)}{\mu_{2}(\lambda)}\right] \cdot \frac{1}{(1-\mu_{1})^{2}}$$

These considerations lead us to the following proposition.

Proposition 4.3.1.

Letting λ^* and $\overline{\lambda}$ be the minimizers of ER(λ) and EV(λ) respectively, the expected inefficiency I^{*} defined by:

$$I^* = \frac{ER(\tilde{\lambda})}{ER(\lambda^*)}$$

satisfies:

$$I^* \leq \frac{1 - h(\lambda^*)}{1 - h(\tilde{\lambda})}$$

and if μ_1 and μ_1^2/μ_2 tend to zero along the sequences $\lambda^*(\nu_0)$ and $\tilde{\lambda}(\nu_0)$ then

 $I^* \rightarrow 1$ as $v_0 \rightarrow \infty$.

Proof:

In view of the approximating assumptions, the proof is identical to Theorem 4.2 of Craven and Wahba's paper. \Box

A NUMERICAL EXPERIMENT

The heuristics are now supplemented by some numerical examples. Two separate logistic regression models were

simulated; independent binomial observations $y_i, y_i \sim B(n, p_i)$ were generated according to:

$$\log \left[\frac{p_{i}}{1 - p_{i}} \right] = x(t_{i}) \quad i = 1, 2, \dots, 80$$
.

The test logit function, x, was:

$$x(t_i) = -3.0 + 4.386 e^{-t_i^2}$$

where $t_i = -1 + \frac{2(i-1)}{79}$, i = 1, 2, ..., 80.

For the first experiment n = 10, for the second n = 1. In each case the logit was estimated as the minimizer in $W_2^2[-1.1]$ of:

$$I_{\lambda}(x) = \sum_{i=1}^{80} \{n \log [1 + e^{x(t_i)}] - y_i x(t_i)\} + \lambda \int_{-1}^{1} [x''(t_i)]^2 dt$$

with λ chosen by the GCV method proposed in section 2. For the logistic model V(λ) reduces to:

$$V(\lambda) = \frac{\sum_{i=1}^{80} \{n\log[1 + e^{x(t_i)}] - y_i x(t_i)\}}{(tr[I - A(\lambda)])^2}$$

where $A(\lambda)$ is the A-matrix (see section 5.2) arising in the final Newton-Raphson minimization problem (see section 3.3), i.e.:

$$\sum_{i=1}^{80} \frac{1}{u_i} [z_i - x(t_i)]^2 + \lambda \int_{-1}^{1} [x''(t)]^2 dt$$

where

$$z_{i} = x_{\lambda}(t_{i}) + \frac{(y_{i} - n\hat{p}_{i})}{n\hat{p}_{i}(1 - \hat{p}_{i})}$$

$$u_{i}^{-1} = n\hat{p}_{i}(1 - \hat{p}_{i})/2$$

and

$$\hat{p}_{i} = \frac{e^{x_{\lambda}(t_{i})}}{1 + e^{x_{\lambda}(t_{i})}}$$

Figure 4.3.1 and 4.3.2 summarize the results of the study. Figure 4.3.1 plots the true and estimated logits while Figure 4.3.2 gives scaled plots of the functions $V(\lambda)$ and $R(\lambda)$. $R(\lambda)$ is given by:

$$R(\lambda) = \sum_{i=1}^{80} np_i (1 - p_i) (x(t_i) - x_{\lambda}(t_i))^2.$$

These plots show that the generalized cross-validation function, $V(\lambda)$, chooses the smoothing parameter quite well in both cases.



Figure 4.3.2

Scaled plots of V and R.



68
CHAPTER 5

A BAYESIAN METHOD FOR INVERTING SATELLITE

RADIANCE DATA

5.1 INTRODUCTION

Remote sensing of the atmosphere is a rapidly developing science. Today's meteorological satellites, such as those in the TIROS-N series, have high resolution instruments on board, which measure the intensity of upwelling radiation at selected channel frequencies. An in-depth description of the data retrieved by the radiometers on the TIROS-N satellites can be found in Smith et. al. (1979). Satellite radiance data make it possible to obtain information on the atmosphere's temperature, moisture, and wind One of the goals of the current Satellite structure. Meteorology programme is to substantially upgrade the quality of the atmospheric information obtained from satellite soundings. A major challenge in this direction is to develop more refined numerical and statistical techniques for inverting the equations of radiative transfer.

For a non-scattering atmosphere in local thermodynamic equilibrium the radiative transfer equations (RTE's) (see Liou 1979) describe how the satellite up-welling radiance measurements relate to the underlying atmospheric tempera-

ture distribution T:

$$R_{v}(T) = B_{v}(T(p_{0}))\tau_{v}(p_{0}) - \int_{0}^{p_{0}} B_{v}[T(p)]\tau_{v}'(p)dp$$

where p_0 is the surface pressure, $\tau_v(p)$ is the transmittance of the atmosphere above pressure p at wavenumber v, and B_v is Plank's function given by:

$$B_{v}[T(p)] = c_{1}v^{3}/[\exp(c_{2}v/T(p)) - 1]$$

$$c_{1} = 1.1906 \times 10^{-5} \text{erg} - \text{cm}^{2} - \text{sec}^{-1}$$

$$c_{2} = 1.43868 \text{cm} - \text{deg}(K)$$

The RTE's are, of course, an idealization. They describe the radiances which the satellite radiometer would record in the absence of such things as atmospheric attenuation due to clouds or instrument noise. Various clever techniques are used to correct for the cloud attenuation problem, after which it is reasonable to model the satellite data as:

 $z_{v_{i}} = R_{v_{i}}(T) + \varepsilon_{i} \quad i = 1, 2, ..., m$

where z_{ν_i} is the satellite "cloud corrected" radiance measurement at wavenumber ν_i , and ε_i is the deviation of z_{ν_i} from the true radiance measurement R_{ν_i} (T). These deviations are probably best thought of as random variables having distributions which can be modeled from practical experience with the radiometer.

A number of procedures have been proposed to retrieve temperature profiles from satellite radiance data (see Fleming (1982) and references cited therein). The Statistical Eigenvector method of Smith and Woolf (1976) is currently favoured. Recently Purser (1983) has pushed for a Bayesian synthesis of meteorological data. In this paper we desribe some Bayesian methods for estimating temperature profiles form the satellite radiance data. Prior information about atmospheric temperature profile structure, channel noise characteristics, and the non-linear dependence of the radiance measurements on temperature are explicitly taken into account. The discussion is restricted to the single profile case but extensions to the estimation of global temperature fields are quite feasible. The basic technique is presented in section 2; in section 3 we examine how a crude version of the Bayesian technique performs on real and simulated satellite radiance data.

5.2 THE BAYESIAN INVERSION TECHNIQUE

5.2.1 GENERAL FORMULATION

Let the temperature profile, T, be written as

 $T = T_0 + \delta$

1) Note added in proof: see also Smith(1983).

where T_0 is an initial guess for T (for example the regional climatology profile) and δ is the update to T_0 which is to be estimated from the data in hand. Our method of estimation consists of considering the functional:

$$I_{\lambda}(\delta) = \frac{1}{m} \sum_{i=1}^{m} w_{i} [z_{i} - R_{v_{i}}(T_{0} + \delta)]^{2} + \lambda J(\delta) \quad \lambda > 0$$

and choosing the update, δ_{λ} , which minimizes I_{λ} over an appropriate space, C, of plausible candidates. The weighted residual sum of squares $\frac{1}{m} \sum_{i=1}^{m} w_i [z_i - R_{v_i}]^2$ is designed to model the information contained in the observed radiance data - there is an implicit assumption that the distribution of the ε_i in (5.1.2) is approximately Gaussian, with mean zero and variance w_i^{-1} . The second term, J, models past information or apriori beliefs about the behaviour of atmospheric temperature profile structure. In a strict Bayesian formulation, J behaves like the -2 log[prior probability density of δ] (see Box and Tiao 1973). The parameter λ specifies the relative weighting given to past and present information. A data-based method for choosing λ is given in section 3.2.

5.2.2 THE CHOICE OF PRIOR

Case A. Smoothness Prior

A general purpose smoothness prior for δ can be

specified by setting

$$J(\delta) = \int_{0}^{P_{0}} [\delta''(p)]^{2} dp.$$

In probabilistic terms (see Kimeldorf and Wahba 1970) this choice of prior says that apriori we believe:

$$(p) = white noise$$
.

The Bayes estimator of δ corresponding to this prior is the minimizer of I_{λ} given by:

$$I_{\lambda}(\delta) = \frac{1}{m} \sum_{i=1}^{m} w_{i} [z_{i} - R_{v_{i}}(T_{0} + \delta)]^{2} + \lambda \int_{0}^{p_{0}} [\delta''(p)]^{2} dp.$$

From a practical point of view, however, the general purpose prior is suboptimal, since a more appropriate representation of prior information would be in terms of regional climatology.

Case B. Extracting a Prior from Regional Climatology

Suppose we represent profiles as a linear combination of basis functions $\{\phi_j\}$. The deviations of a profile from regional climatology can then be written as:

$$\delta = \sum_{j} \alpha_{j} \phi_{j}$$

Regional climatology data bases can now be used to model the distribution of the basis coefficient vector, α . If

the number of basis functions is finite, and the distribution of α was modelled by a multivariate normal distribution

$$\alpha = N(0, \Sigma_{\alpha})$$

where Σ_{α} is the sample covariance matrix obtained from the regional climatology data base, the resulting Bayesian estimator of δ would be:

$$\delta_{\lambda} = \sum_{j} \hat{\alpha}_{j} \phi_{j}$$

where $\hat{\alpha}$ is the minimizer of:

$$I_{\lambda}(\alpha) = \frac{1}{m} \sum_{i=1}^{m} w_{i} [z_{i} - R_{v_{i}} (T_{0} + \sum_{j} \alpha_{j} \phi_{j})]^{2} + \lambda \alpha^{t} \Sigma_{\alpha}^{-1} \alpha.$$

Under either of the above two priors, the iterative scheme used to approximate δ_{λ} is essentially the same. At each stage the non-linear radiance functionals, R_{ν_i} ($T_0 + \delta$), are linearized about the current iterate, δ^k , and the resulting quadratic minimization problem is solved.

5.3 COMPUTATIONAL METHODS

5.3.1 THE GAUSS-NEWTON ALGORITHM

An iterative Gauss-Newton procedure can be used to obtain the optimal correction defined by the minimizer of I_1 . Let T^k be the temperature profile estimate at the kth iteration.

 $\mathbf{T}^{\mathbf{k}} = \mathbf{T}_{\mathbf{0}} + \mathbf{\delta}^{\mathbf{k}}$

The iterative algorithm for obtaining the Bayes estimate can be stated as follows:

Iterative Gauss-Newton Procedure while [convergence criteria fail to be met] do

Linearize the radiance functionals Solve for δ^{k+1}

Update convergence indicators

Two cases corresponding to the smoothnesss [and regional

end

F

A detailed description of the iterative steps follows.

STEP 1: LINEARIZATION OF THE RADIANCE FUNCTIONALS

Linearizing R_{v_i} (T) about T^k we obtain:

 $R_{v_{i}}(T) \approx R_{v_{i}}(T^{k}) + R_{v_{i}}(T_{k}) (\delta - \delta^{k})$

where the action of the linear functional $R_{\nu_i}^{\prime}(T^k)$ on a function h is given by:

$$R_{v_{i}}^{\prime}(T^{k})h = K_{s}(v_{i}, p_{0})h(p_{0}) - \int_{0}^{p_{0}} K(v_{i}, p)h(p)dp$$

The kernels K_s and K are:

$$K_{s}(v_{i}, p_{0}) = B_{v_{i}}'[T^{k}(p)]\tau_{v_{i}}(p_{0})$$
$$K(v_{i}, p) = B_{v_{i}}'[T^{k}(p)]\tau_{v_{i}}(p_{0})$$

STEP 2: SOLVING FOR THE NEXT ITERATE

After the linearization we get the next iterate, $\delta^{k+1},$ by minimizing the functional

$$I_{\lambda}^{(k)}(\delta) = \sum_{i=1}^{m} w_{i} [z_{i}^{k} - R_{\nu_{i}}'(T^{k})\delta]^{2} + \lambda J(\delta)$$
$$z_{i}^{k} = z_{i} - R_{\nu_{i}}(T^{k}) + R_{\nu_{i}}'(T^{k})\delta^{k} \quad i = 1, 2, ..., m$$

Two cases corresponding to the smoothnesss and regional climatology priors will be considered.

Solution for the Smoothness Prior:

Here δ^{k+1} is the minimizer in $C = \{\delta: \int_0^{p_0} [\delta''(p)]^2 dp < \infty\}$ of the functional:

$$I_{\lambda}^{(k)}(\delta) = \frac{1}{m} \sum_{i=1}^{m} w_{i} [z_{i}^{k} - R_{v_{i}}^{i}(T^{k})\delta]^{2} + \lambda \int_{0}^{p_{0}} [\delta^{"}(p)]^{2} dp.$$

It follows by straightforward calculus that δ^{k+1} can be written as:

$$\delta^{k+1}(p) = \hat{d}_1 + \hat{d}_2 p + \sum_{j=1}^{m} \hat{c}_j \xi_j(p)$$

where $\xi_j(p) = R_{v_i}'(T_k)Q_p(\cdot)$ and

$$Q_{p}(r) = \begin{cases} p^{2}r/2 - r^{3}/6, & p \ge r \\ \\ pr^{2}/2 - p^{3}/6, & p \le r \end{cases}$$

Letting $K_{ij} = \langle \xi_i, \xi_j \rangle = \int_0^{p_0} \xi_i^* \xi_j^* dp$ for i, j = 1, 2, ..., m, and $T_{1i} = 1$, $T_{2i} = R'_{v_i} (T^k) p$ for i = 1, 2, ..., m, the coefficients \hat{d} and \hat{c} are the solution to the linear system:

$$[K + m\lambda I]_{C} + Td = z^{K}$$
$$T^{t}_{C} = 0$$

Taking a QR-decomposition of the T matrix $(T = [Q_1:Q_2])$ R with $R^t = [R_1:0]$ see Bunch et. al (1979)) \hat{d} and \hat{c} are given by the formulae:

$$\hat{c} = Q_2^{t} [Q_2^{t} K Q_2 + m\lambda I]^{-1} Q_2 z^{k}$$
$$\hat{d} = R_1^{-1} Q_1^{t} [z^{k} - K\hat{c}]$$

Note that the calculation of K_{ij} can be done by first fitting cubic spline interpolants, $\tilde{\xi}_{j}$, to ξ_{j} evaluated on a fine grid and proceeding to approximate K_{ij} as:

$$K_{ij} \approx \langle \tilde{\xi}_{i}, \tilde{\xi}_{j} \rangle$$

This method of approximation is used by Nychka (1983).

Solution for the Regional Climatology Prior:

 δ^{k+1} is now given by

$$\delta^{k+1} = \sum_{j} \alpha_{j}^{k+1} \phi_{j}$$

where α^{k+1} is the minimizer of:

$$I_{\lambda}^{k}(\alpha) = \frac{1}{m} \sum_{i=1}^{m} w_{i} [z_{i}^{i} - X_{i}\alpha]^{2} + \lambda \alpha^{t} \Sigma_{\alpha}^{-1} \alpha.$$

So $\alpha^{k+1} = [X^{t}D_{\omega}X + m\lambda \Sigma_{\alpha}^{-1}]^{-1}X^{t}D_{\omega} z^{k}$ where X and D_{ω}
are:

$$\begin{aligned} \mathbf{x}_{ij} &= \mathbf{K}_{s} (\mathbf{v}_{i}, \mathbf{p}_{0}) \phi_{j} (\mathbf{p}_{0}) - \int_{0}^{\mathbf{p}_{0}} \mathbf{K} (\mathbf{v}_{i}, \mathbf{p}) \phi_{j} (\mathbf{p}) d\mathbf{p} \\ \mathbf{D}_{\omega} &= \operatorname{diag} (\mathbf{w}_{1}, \mathbf{w}_{2}, \dots, \mathbf{w}_{m}) . \end{aligned}$$

The kernels K_s and K are given above.

STEP 3: CONVERGENCE CRITERIA

Standard convergence criteria can be used in the above iterations. We have used:

$$e_{1}^{k+1} = |I_{\lambda}(\delta^{k+1}) - I_{\lambda}(\delta^{k})| / I_{\lambda}(\delta^{k})$$
$$d_{2}^{k+1} = ||\delta^{k+1} - \delta^{k}|| / ||\delta^{k}||$$

where $\|\cdot\|$ is an appropriate norm.

5.3.2 CHOICE OF THE SMOOTHING PARAMETER

An extension of Generalized Cross Validation (GCV) can be used to estimate the smoothing parameter, λ , empirically from the data. GCV was developed by Craven and Wahba (1979) to estimate a smoothing parameter similar to λ in a linear functional model context. In their model, the observations z_i are noise contaminated linear functionals of an unknown function, x, of interest.

 $z_{i} = L_{i}x + noise, i = 1, 2, ..., m$

where the L_i 's are linear functionals. The estimator, x_{λ} , of x is the minimizer of the functional

 $\frac{1}{m} \sum_{i=1}^{m} [z_i - L_i x]^2 + \lambda \int [x''(t)]^2 dt \quad \lambda > 0.$

It turns out that the residuals, $z_i - \hat{z}_i$, can be written as a linear transform of the data.

 $z - \hat{z} = [I - A(\lambda)]z$

where $A(\lambda)$ is an mxm matrix defined in terms of the linear functionals. Formally $A(\lambda)_{ij} = \frac{\partial L_i}{\partial z_j} (x_{\lambda})$. The GCV estimate of λ is the minimizer of the function

$$V(\lambda) = \frac{\frac{1}{m} \sum_{i=1}^{n} [z_i - \hat{z}_i]^2}{(\frac{1}{m} \operatorname{tr}[I - A(\lambda)])^2}$$

tone continuing of

$$\frac{\frac{1}{m} \operatorname{RSS}(\lambda)}{\left[1-p(\lambda)\right]^2}$$

 $V(\lambda)$ is an assessment of λ combining goodness of fit and model dimensionality (Wahba refers to tr[I-A(λ)] as the "equivalent degrees of freedom" of RSS(λ)) of model. The goodness of fit is measured by the residual sum of squares, $\frac{1}{m}$ RSS(λ), while the dimensionality of x_{λ} is measured by the "effective" number of parameters describing x_{λ} relative to the sample size, i.e. $p(\lambda)$. As $\lambda \rightarrow \infty$, x_{λ} tends to a straight line regression making $\frac{1}{m}$ RSS(λ) take on its maximum and $p(\lambda)$ be a minimum. On the other hand, as $\lambda \rightarrow 0$, x_{λ} tends to a data interpolant minimizing $\frac{1}{m}$ RSS(λ) and maximizing $p(\lambda)$. It has been shown that the GCV estimate of λ , for large sample sizes, comes close to minimizing the integrated squared error between x and x_{λ} :

$$R(\lambda) = \frac{1}{m} \sum_{i=1}^{m} [L_i x - L_i x_{\lambda}]^2 .$$

For further discussion see Craven and Wahba (1979), Speckman (1983) and Wahba (1982).

Extended GCV Estimator of the Smoothing Parameter

In the linear functional model the definition of

dimensionality is facilitated by linearity. Even though the radiance functionals are not linear a simple approximation can be used to develop a dimensionality measure similar in spirit to $p(\lambda)$ above. This approximation makes use of the fact that the optimal correction under the smoothness and regional climatology priors minimize quadratic functionals:

$$\frac{1}{m} \sum_{i=1}^{m} w_{i} [z_{i} - R_{v_{i}}(T_{\lambda}) + R_{v_{i}}'(T_{\lambda})\delta_{\lambda} - R_{v_{i}}'(T_{\lambda})\delta]^{2} + \lambda J(\delta) \qquad (*)$$

so that the residuals are defined, implicitly, by a linear equation.

$$z - R_{1}(T_{1}) = [I - A_{1}(\lambda)]z^{*}$$

where $z_i^* = z_i - R_{v_i}(T_{\lambda}) + R_{v_i}(T_{\lambda})\delta_{\lambda}$ and $A_1(\lambda)$ is the A-matrix arising in the minimization of (*).

We therefore have, that to a first order approximation, the effective number of parameters in δ_{λ} relative to the sample size is given by:

$$p_1(\lambda) = tr[A_1(\lambda)]/m$$

which leads us to the GCV-type assessment of λ as

$$V(\lambda) = \frac{\frac{1}{m} RSS(\lambda)}{[1-p_1(\lambda)]^2}$$

with the KIRS (the

Note that the matrix $A_1(\lambda)$ is already available from the final stage of the Gauss-Newton iteration. After some algebra one can show that $V(\lambda)$ for the smoothness prior is given by:

$$V(\lambda) = \frac{\|\hat{c}\|_{R^{m}}^{2}}{\frac{1}{m} [tr(Q_{2}^{t}KQ_{2} + m\lambda I)^{-1}]^{2}}$$

and for the regional climatology prior $V(\lambda)$ becomes:

$$V(\lambda) = \frac{\frac{1}{m} \sum_{i=1}^{m} w_i [z_i - R_{v_i}(T_{\lambda})]^2}{[\frac{1}{m} \operatorname{tr} \Sigma_{\alpha}^{-1} [X^{\dagger} D_{\omega} X + m \lambda \Sigma_{\alpha}^{-1}]^{-1}]^2}$$

5.4 PERFORMANCE OF THE TECHNIQUE

In order to assess the potential of the Bayesian method described in section 2, analyses with real and simulated radiance measurements were performed. For each of 15 locations shown in Figure 5.4.1, an operational profile (obtained via the statistical eigenvector method of Smith and Woolf 1976) and a verification profile (obtained from radiosonde data) were kindly made available to us by Dr. Tom Koehler of the Meteorology Department at UW-Madison. Dr. Koehler also provided us with the HIRS (the seven 15µm and five 4.3µm band channels) and MSU (the three O₂ channels) satellite



0027 January 1980 TIROS-N Sounding Locations



 \circ clear, \triangle partly cloudy, + cloudy.

The 15 shaded circles were the locations used in the study.

radiance data corresponding to the operational profile. These soundings were given "cloud free" classification by the HIRS radiometer.

The radiance data appeared to have systematic biases by channel. Figure 5.4.2 plots the equivalent brightness temperature differences between the observed radiance measurements and their "true" values, i.e. the radiances predicted by the verification profile. (The brightness temperature of a radiance measurement z at wavenumber v is defined to be $B_{y}^{-1}(z)$ i.e. the equivalent black-body temperature corresponding to z.) Table 5.4.1 gives the median bias and the scale of measurement error computed from Figure 5.4.2. Using the table below, a simple correction was made to the radiance data - a raw radiance measurement was corrected by subtracting off the median bias found in the channel. In addition to these bias corrected radiance data, simulated radiance measurements were generated from the verification profile at each location according to:

$$z_i = R_{v_i}(T) + \sqrt{w_i} \varepsilon_i$$

where $\varepsilon_i \sim N(0,1)$ and the $\sqrt{w_i}$'s corresponded to the standard channel weights given in column 4 of Table 5.4.1 i.e. column 4 converted into radiance units. T is the verification profile.

Brightness temperature differences between satellite radiance measurements and "true" radiance i.e. those predicted by the verification profiles.



TABLE 5:4.1

Channel Number from Smith et al (1979)	Channel Wave Number	Estimated Bias from Figure l	Estimated Noise level from Figure l	Noise Level from Past Experience
oineatilite	7 edd od	HIRS Chan	nels	
s del esw s	668	-0.62	1.54	2.50
2	679	-1.38	1.32	1.20
	691	-0.68	0.68	1.30
4	704	-0.38	0.38	0.75
5	716	-1.02	0.47	1.00
6	732	-1.34	0.66	1.00
boog 7 as a	748	-2.25	1.49	1.50
13	2190	-0.40	1.74	1.00
14	2213	-0.30	1.14	0.70
15	2240	0.08	0.82	1.00
16	2276	3.16	3.46	1.40
17	2361	4.07	4.34	3.00
		MSU Chann	lels	Nerro ensi
2	1.792	0.11	1.12	0.60
3	1.833	0.52	1.01	0.50
4	1.933	-0.84	0.50	1.00

SUMMARY STATISTICS FROM FIGURE 5.4.2

Using the algorithm in section 2, with the general purpose prior and λ chosen by the modified GCV criterion, two temperature profiles were estimated at each location - one from bias corrected satellite data, and the other from simulated radiances. Throughout these calculations transmittances corresponding to the verification profile were used, and the surface pressure was set at 1000mb. The operational sounding was used as the first guess.

The individual retrievals are given in an appendix at the end of this chapter. Figures 5.4.3 - 5.4.5 summarize the results of the study; three average bias and variability plots are presented. Bias and variability were computed as follows: for the ith map location, let δ_i be the difference between the verification profile and the retrieved profile. The bias, b(p), at pressure p is the 20% trimmed mean (mean taken over the middle 60% of values) of $\{\delta_i(p)\}_{i=1}^{15}$, while the variability is the 20% trimmed mean of the absolute deviations $\{|\delta_i(p) - b(p)|\}_{i=1}^{15}$. Because of the presence of a few "wild" δ_i 's it was felt that these were more accurate estimates of bias and variability than the usual mean and standard deviation (see Tukey 1977).

Figure 5.4.3

Performance of the Bayesian Method

on Simulated Radiance Data



Figure 5.4.4

Performance of the Operational Method on Actual Satellite Radiance Data



Figure 5.4.5

Performance of the Bayesian Method on Bias-corrected Satellite Radiance Data



The Bayesian retrievals perform quite well, especially since no real regional climatology has been incorporated into the prior specification. While the operational retrieval has a large bias in the neighbourhood of the tropopause, the Bayesian retrievals are much less biased at this level. At other levels the Bayesian method (at least with simulated noisy radiance data) does just as well as the operational retrieval. The level of variability in all three retrieval methods is approximately 2°K. The variability in the Bayesian retrieval with satellite (bias corrected) radiances has a peak near the surface. This peak may be due, in part, to the fact that we assumed surface pressure was 1000mb which for retrievals near the Rockies is somewhat suspect.

Overall the plots are encouraging and suggest that it may be worthwhile investing some more time into developing a Bayesian retrieval method using a prior derived from regional climatology.

APPENDIX B. RESULTS OF INDIVIDUAL RETRIEVALS

Fifteen figures follow. Each figure has two plots corresponding to the retrievals from simulated radiance data and bias-corrected satellite radiance data respectively. The true profile (solid), the operational profile (dot), and the estimated profile (dash) are all given.

The latitude (degrees North) and longitude (degrees East) corresponding to each retrieval are as follows:

Retrieval #	Latitude	Longitude
1	30.48	95.36
2	31.28	92.51
3	31.68	89.95
4	31.99	85.99
5	32:83	95.41
6	33.37	94.08
7	34.41	85.15
8	37.73	105.27
9 .	39.52	98.05
10	39.95	96.15
11	39.04	106.74
12	39.91	99.51
13	47.39	109.69
14	48.73	110.28
15	51.36	101.99







Figure 5.A.2: Results for Retrieval 2



Figure 5.A.3: Results for Retrieval 3



Figure 5.A.4: Results for Retrieval 4















Figure 5.A.8: Results for Retrieval 8



Figure 5.A.9: Results for Retrieval 9



Figure 5.A.10: Results for Retrieval 10














Figure 5.A.14: Results for Retrieval 14



Figure 5.A.15: Results for Retrieval 15

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APPENDIX C

EXISTENCE OF SMOOTH SOLUTIONS TO THE RADIATIVE TRANSFER EQUATIONS

Here we show how the non-linear regression theory developed in section 3 of Chapter 2 can be used to establish the existence of smooth (smoothness prior) solutions to the ill posed problem associated with the radiative transfer equations. Recall that the objective functional here has the form

$$I_{\lambda}(\delta) = \sum_{i=1}^{n} [z_{i} - R_{v_{i}}(T_{0} + \delta)]^{2} + \lambda \int_{0}^{P_{0}} [\delta^{(m)}(p)]^{2} dp \qquad (B.1)$$

and we are looking for minimizer of this functional in C_{T_0} where $C_{T_0} = \{\delta \in W_2^m[0,p_0] | T_0 + \delta \ge 0\}; T_0$ is the initial guess profile and $R_v(T)$ is given by:

$$R_{v}(T) = B_{v}[T(p_{0})]\tau_{v}(p_{0}) - \int_{0}^{p_{0}} B_{v}[T(p)]T_{v}'(p)dp$$

where B_v is Plank's function and τ_v is atmospheric transmittance as a function of pressure.

We note from physical considerations that atmospheric transmittance is a continuously differentiable and monotonic decreasing function of pressure. We assume that the initial guess profile is continuous in pressure. These

109

facts will be made use of in the sequel.

For each $m \ge 1$, C_{T_0} is clearly a convex set but since evaluation is a continuous linear functional in W_2^m for $m \ge 1$ it follows that C_{T_0} is also weakly closed. Let us now make the following three claims:

- (c.1) R_{v} is weakly continuous on $C_{T_{0}} \forall v > 0$
- (c.2) $R_{v}(T_{0} + \delta) = 0 \quad \forall \quad \delta \in C_{T_{0}}$ with equality iff $T_{0} + \delta \equiv 0$ (c.3) $R_{v}(T_{0} + \delta)$ is convex as a functional of δ on C_{m}

Our existence result will follow from these claims.

Theorem B.l.

If (c.l)-(c.3) are true then $\exists \delta_{\lambda} \in C_{T_{0}}$ such that

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$$I_{\lambda}(\delta_{\lambda}) = \min_{\substack{\delta \in C_{T_{0}}}} I_{\lambda}(\delta)$$

Proof:

From (c.1) we have by Theorem 2.2 that I_{λ} is weakly lower semi-continuous on C_{T_0} also using the continuity of evaluation C_{T_0} satisfies property 2 of section 2 with $z \equiv -T_0$. Letting $\phi(x) = |x|$ in Theorem 2.3.5 we have from (c.2) and (c.3) that $\sum_{i=1}^{n} \phi[R_{\nu_i}(T_0 + \delta)]$ is convex and $\sum_{i=1}^{n} \phi[R_{\nu_i}(T_0 + \delta)] = 0 \iff T_0 + \delta = 0$ (i.e. $\sum_{i=1}^{n} R_{\nu_i}(T_0 + \delta)$ has a proper minimizer, $-T_0$, in $[-T_0 + H_0] \cap C_{T_0}$. Hence I_{λ} is coercive on C_{T_0} and the result follows by Theorem 2.2.1. \Box

The remainder will be devoted to proving (c.1)-(c.3). We begin with a few elementary lemmas concerning the Plank's function B_{y} .

Lemma B.2.

 $\forall v > 0 \quad B_v$ is a convex function on $[0, \infty)$.

Proof:

$$B_{\nu}(t) = c_1 v^3 / [exp(c_2 v/t) - 1]$$

Clearly it is enough to show that

f(x) = 1/[exp(1/x) - 1] is convex for $x \ge 0$.

A direct calculation reveals that

$$f''(x) = \frac{\exp(1/x)}{x^{3} [\exp(1/x) - 1]^{2}} g(1/x)$$

where

$$g(y) = -2 - y + 2y \exp(y) / [\exp(y) - 1]$$
$$= [(y - 2)(\exp(y) - 1) + 2y] / [\exp(y) - 1]$$

but $1/[exp(y) - 1] > 0 \quad \forall y > 0$ and

$$(y-2) [\exp(y) - 1] + 2y = \sum_{j=3}^{\infty} \frac{y^{j}[j! - 2(j-1)!]}{j!(j-1)!}$$
$$= \sum_{j=3}^{\infty} \frac{y^{j}(j-2)}{j!}$$
$$\ge 0 \ \forall y \ge 0 \ .$$

Consequently $g(1/x) \ge 0 \forall x \ge 0$ and this implies that $f''(x) \ge 0 \forall x \ge 0$. It follows that f and therefore B_v are convex on $[0, \infty)$.

Lemma B.3.

 B_{v}' is bounded on $[0,\infty)$.

Proof:

$$B_{v}'(t) = \frac{c_{1}c_{2}v^{4}\exp(c_{2}v/t)}{t^{2}\left[\exp(c_{2}v/t) - 1\right]^{2}}.$$

This is continuous in t and tends to $c_1^{\nu^2/c_2}$ as $t \neq \infty$. Consequently B'_{ν} is bounded on $[0,\infty)$.

Lemma B.4.

[Justification of (c.l)] $\forall v > 0$, $R_v(T_0 + \delta)$ is weakly continuous on C_{T_o} .

Proof:

$$R_{v}(T) = B_{v}[T(p_{0})]\tau_{v}(p_{0}) - \int_{0}^{p_{0}} B_{v}[T(p)]\tau_{v}'(p)dp$$

Now evaluation is a continuous linear functional on $W_2^{m}[0,p_0]$ and B_v is continuous on $[0,\infty)$ it follows that $B_v[T(p_0)]\tau_v(p_0)$ is weakly continuous on C_{T_0} . It remains to show that $\int_0^{p_0} B_v[T(p)]\tau'_v(p)dp$ is weakly continuous. Let $\{\delta_n\}$ be a sequence of functions in C_{T_0} converging weakly to δ (δ must also lie in C_{T_0} by weak closedness) and let $T_n = T_0 + \delta_n$ and $T_0 + \delta$.

$$\begin{aligned} & \left| \int_{0}^{P_{0}} B_{v}[T_{n}(p)]\tau_{v}'(p)dp - \int_{0}^{P_{0}} B_{v}[T(p)]\tau_{v}'(p)dp \right| \\ & \leq \int_{0}^{P_{0}} \left| B_{v}[T_{n}(p)] - B_{v}[T(p)] \| \tau_{v}(p) \right| dp \\ & \leq \text{ constant } \int_{0}^{P_{0}} \left| T_{n}(p) - T(p) \| \tau_{v}'(p) \right| dp \end{aligned}$$

by Taylor's theorem and the boundedness of B'_{v} (Lemma B.3). τ_{v} is a transmittance function so τ'_{v} is continuous which implies that τ'_{v} is bounded on $[0,p_{0}]$. Consequently

$$\int_{0}^{p_{0}} |\mathbf{T}_{n}(\mathbf{p}) - \mathbf{T}(\mathbf{p})| |\tau_{v}'(\mathbf{p})| d\mathbf{p} \leq \text{constant } \int_{0}^{p_{0}} |\mathbf{T}_{n}(\mathbf{p}) - \mathbf{T}(\mathbf{p})| d\mathbf{p}$$

Hence

$$\int_{0}^{p_{0}} B_{v}[T_{n}(p)]\tau_{v}'(p)dp \rightarrow \int_{0}^{p_{0}} B_{v}[T(p)]\tau_{v}'(p)dp \text{ as } n \rightarrow \infty$$

i.e. $\int_0^{p_0} B_v[T(p)]\tau'_v(p)dp$ is weakly continuous on C_{T_0} . Finally since the sum of two weakly continuous functionals is weakly continuous it follows that R_v is weakly continuous on C_{T_0} .

Lemma B.5.

[Justification of (c.2)] $\forall v > 0$, $R_v(T_0 + \delta) \ge 0$ $\forall \delta \in C_{T_0}$ with equality iff $T_0 + \delta = 0$.

Proof:

Positivity follows from the positivity of B_v on $[0,\infty)$ and the fact that τ_v is a positive strictly monotonic decreasing function of pressure. Clearly $R_v(0) = 0$ since $B_v(0) = 0$. Now suppose $R_v(T_0 + \delta) = 0$ for some $\delta \in C_{T_0}$. Then $B_v[T_0(p_0) + \delta(p_0)]\tau_v(p_0) = 0 \iff T_0(p_0) + \delta(p_0) = 0$ and $\int_0^{p_0} B_v[T_0(p) + \delta(p)]\tau_v'(p)dp = 0$. If $\Xi p^* \in [0, p_0)$ such that $T_0(p^*) + \delta(p^*) > 0$ then by continuity of $T_0 + \delta \equiv \varepsilon > 0$ such that $\forall p \in [p^*, p^* + \varepsilon]$ $T_0(p) + \delta(p) > 0$. But, since B_v and $-\tau_v'$ are both strictly positive functions on their respective ranges, we would then have that

$$- \int_{0}^{p_{0}} B_{v} [T_{0}(p) + \delta(p)] \tau_{v}'(p) dp > 0 \implies R_{v} (T_{0} + \delta) > 0 .$$

This is a contradiction. Consequently $R_{\nu}(T_0 + \delta) = 0$ iff $T_0 + \delta = 0$.

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ABSTRACT

There are many areas in applied science where the non-parametric estimation of regression functions is important.

In this thesis a general penalized likelihood method for non-parametrically estimating regression functions under a variety of observational models is developed. The existence and numerical approximation of the estimators is studied and a cross-validatory method for estimating the smoothing parameter is presented. Implementation of the method is algorithmically straightforward.

The procedures developed are applied to the estimation of atmospheric temperature profiles from satellite radiance data and are found to compare favorably with the currently used methodology.