# Efficient Estimation of a Semiparametric <br> Characteristic-Based Factor Model of Security Returns 

By<br>Gregory Connor<br>Matthias Hagmann<br>Oliver Linton

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Gregory Connor is Professor of Finance in the Department of Finance at the London School of Economics and Political Science (LSE), where he has been teaching since 1991. He is also Director of the Asset Pricing and Portfolio Management research programme for the Financial Markets Group. Gregory served as full-time Director of Research-Europe in Barra's London office from 1994 until 1999, while keeping his LSE link as visiting lecturer. He returned to LSE full-time in 1999. Prior to joining LSE, Gregory was on the finance faculty at Northwestern University and the University of California Berkeley. His research focuses on security market risk and return, and related portfolio management topics. After graduating from the MSc Econometrics and Mathematical Economics (with distinction) at the London School of Economics and Political Science, Matthias Hagmann enrolled into the Finance PhD Program at the Swiss Finance Inistitute. He graduated at the University of Geneva in November 2007 under the supervision of Prof Olivier Scaillet. Matthias' academic work has been published so far in the Journal of Econometrics and the Journal of Business and Economic Statistics. He works full time as a Quantitative Analyst for Concordia Advisors, a 2bn \$ multi-strategy hedge fund with offices in London, New York, Singapore and Bermudas. Oliver Linton is Professor of Econometrics in the Department of Economics, London School of Economics. He is Co-editor of Econometric Theory and Associate Editor of: Econometrica, The Journal of Econometrics and of The Journal of the American Statistical Association. Any opinions expressed here are those of the authors and not necessarily those of the FMG. The research findings reported in this paper are the result of the independent research of the authors and do not necessarily reflect the views of the LSE.

# Efficient Estimation of a Semiparametric Characteristic-Based Factor Model of Security Returns* 

Gregory Connor<br>The London School of Economics<br>Matthias Hagmann<br>Concordia Advisors and Swiss Finance Institute<br>Oliver Linton ${ }^{\dagger}$<br>The London School of Economics

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#### Abstract

This paper develops a new estimation procedure for characteristic-based factor models of security returns. We treat the factor model as a weighted additive nonparametric regression model, with the factor returns serving as time-varying weights, and a set of univariate nonparametric functions relating security characteristic to the associated factor betas. We use a time-series and cross-sectional pooled weighted additive nonparametric regression methodology to simultaneously estimate the factor returns and characteristic-beta functions. By avoiding the curse of dimensionality our methodology allows for a larger number of factors than existing semiparametric methods. We apply the technique to the three-factor Fama-French model, Carhart's four-factor extension of it adding a momentum factor, and a five-factor extension


[^0]adding an own-volatility factor. We find that momentum and own-volatility factors are at least as important if not more important than size and value in explaining equity return comovements. We test the multifactor beta pricing theory against the Capital Asset Pricing model using a standard test, and against a general alternative using a new nonparametric test.

Keywords: Additive Models; Arbitrage pricing theory; Factor model; Fama-French; Kernel estimation; Nonparametric regression; Panel data.

JEL CODES: G12, C14.

## 1 Introduction

This paper has two objectives. First, to develop practically relevant econometric methodology for estimation of a semiparametric panel data model in the case where both cross-section and time series dimensions are large. Second, to use this methodology to address some central questions in empirical finance.

Individual stock returns have strong common movements, and these common movements can be related to individual security characteristics such as market capitalization and book-to-price ratios. Rosenberg (1974) develops a factor model of stock returns in which the factor betas of stocks are linear functions of observable security characteristics. Rosenberg's approach requires a strong assumption of linearity, but is thenceforward a standard econometric model. Fama and French (1993) use portfolio grouping to estimate a characteristic-based factor model without assuming linearity. They estimate a three-factor model, with a market factor, size factor and value factor. The market factor return is proxied by the excess return to a value-weighted market index. The size factor return is proxied by the difference in return between a portfolio of low-capitalization stocks and a portfolio of highcapitalization stocks, adjusted to have roughly equal book-to-price ratios. The value factor is proxied by the difference in return between a portfolio of high book-to-price stocks and a portfolio of low book-to-price stocks, adjusted to have roughly equal capitalization. Using these factor returns, the factor betas are estimated via time-series regression. Their methodology is very influential in empirical finance, their papers being amongst the most cited finance papers of all time. However, there does not seem to be a way of obtaining consistent standard errors in their framework.

Connor and Linton (2007) develop a semiparametric methodology that combines elements of the Rosenberg and Fama-French approaches. They describe a characteristic-based factor model like Rosenberg's but replacing Rosenberg's assumption that factor betas are linear in the characteristics with an assumption that factor betas are smooth nonlinear functions of the characteristics. In a model with two characteristics, size and value, plus a market factor, they form a grid of equallyspaced characteristic-pairs. They use multivariate kernel methods to form factor-mimicking portfolios for the characteristic-pairs from each point on the grid. Then they estimate factor returns and factor betas simultaneously using bilinear regression applied to the set of factor mimicking portfolio returns.

A weakness of the Connor-Linton methodology is the reliance on multivariate kernel methods to create factor-mimicking portfolios. These multivariate kernel methods severely restrict the number of factors which can be estimated well using their technique due to the curse of dimensionality, see Stone (1980). The same problem appears in a different guise with the Fama-French methodology. To create their size and value factor returns, Fama and French double-sort assets into size and value categories. Adding a third characteristic with this method requires triple-sorting, adding a fourth
requires quadruple-sorting; like Connor-Linton, the method quickly becomes unreliable for typical sample sizes and more than two characteristic-based factors.

In this paper we develop a new estimation methodology that does not require any portfolio grouping or multivariate kernels. Instead, we estimate the factor returns and characteristic-beta functions using the methodology of weighted additive nonparametric regression. This relies on the fact that in each time period the characteristic-based factor model proposed in Connor-Linton is a weighted additive sum of univariate characteristic-based functions. The nonparametric part of the estimation problem is made univariate by decomposing the full problem into an iterative set of sub-problems in each characteristic singly, a standard trick in weighted additive nonparametric regression. We modify the weighted additive nonparametric regression methodology to account for our model's feature that the weights vary each time period while the characteristic-beta functions stay constant. Our estimation approach is to derive population moment conditions for the unknown quantities and then solve them with empirical counterparts. The moment conditions we obtain for the characteristic functions constitute a system of linear type II Fredholm integral equations, and so falls in the class of well-posed inverse problems reviewed in Carrasco, Florens, and Renault (2006, section 7). Recently, Mammen, Linton, and Nielsen (1999) have provided a pointwise distribution theory for a special case of this problem for a single cross-section. What is new in this paper is that we have a panel setting with a large cross-section and time series and our model allows for an increasing number of "parametric" quantities, i.e., the unobserved factors that have to be estimated at the same time.

Our model falls in the class of semiparametric panel data models for large cross-section and long time series. There has been some work on semiparametric models for panel data, see for example Kyriazidou (1997), nonparametric models, see for example Bester and Hansen (2007), and nonparametric additive models, see for example Porter (1996) and more recently Mammen, Støve, and Tjøstheim (2006). Most of this work is in the context of short time series, and with a consequent focus on microeconometric issues. More recently, there has been work on panel data with large crosssection and time series dimension, especially in finance where the datasets can be large along both dimensions and in macroeconomics where there are cross-sectional panels of many related series (such as business conditions survey data) with quite long time series length. Phillips and Moon (1999) investigated the consequence of having a large cross-section for the time series issue of nonstationarity; they develop probabilistic results for such settings and establish limit theorems for linear regression estimators. In a series of papers, Bai and $\operatorname{Ng}(2002)$, Bai $(2003,2004)$ have recently developed, under very weak conditions, an extensive inferential theory for a class of statistical factor models that are widely used in empirical finance. Our model differs from theirs in that we include observable covariates, which changes some identification issues. The estimation methods and asymptotic results
are also quite different. Pesaran (2006) proposes a new method for estimation of linear panel data models with observed and unobserved covariates in this setting using cross-sectional averaging. Our theoretical analysis uses the "transpose" of his method, i.e., time series averaging, to obtain simple initial consistent estimators for our model. We work with a semiparametric model that allows more flexibility with regard to the functional form of the observable covariate effect. This can be important in practice. We establish a number of statistical results for our method that can be used to conduct inference or aid implementation. We establish the pointwise asymptotic normality of the functional components of our model at what appears to be an optimal rate. We further show that our method has a sort of oracle efficiency property. We also establish the asymptotic normality of the estimated factors. The small sample properties of related procedures have been thoroughly analyzed by Monte Carlo methods, see for example Hastie and Tibshirani (1990), Härdle, Linton, and Sperlich (1999), and Linton and Mammen (2005).

Our model allows for any number of factors with no theoretical loss of efficiency, and we exploit this in our application. In addition to the market, size and value factors of the standard Fama-French model, we add a momentum factor as suggested by Jagadeesh and Titman (1993) and Carhart (1997), and an own-volatility factor, a choice influenced by the recent work of Goyal and Santa Clara (2003) and Ang, Hodrick, Xing and Zhang (2006a, 2006b). This reflects the feature that our methodology allows us to estimate a model with more factors. We find that the two added factors, momentum and volatility, are as important or more important than size and value in explaining equity return comovements. Hence, the improved data-efficiency of our new method has real empirical value.

We test the single-factor Capital Asset Pricing Model against a five-factor asset pricing model using the standard Gibbons, Ross and Shanken (1989) test; we find that the CAPM is rejected. We also develop a new nonparametric test for multifactor pricing models as part of our estimation methodology. To implement the test we assume that mispricing is a smooth multivariate function of observable security characteristics. We estimate this mis-pricing function simultaneously with the factor model of returns, and test whether the mispricing function is the null function. We find that the five-factor model does a good job of explaining asset return premia; the alpha function differs only negligibly from a null function, at least for the four security characteristics that we consider.

Finally, we evaluate various time series models for the risk factors. We establish the asymptotic properties of two-stage estimators of the parameters of this model and estimate vector autoregressions both for the levels of the factor returns and the factor returns squared to explore factor return and volatility dynamics.

We proceed as follows: Section 2 presents the model. Section 3 describes the estimation algorithm in the balanced and unbalanced panel case. Section 4 develops the distribution theory. Section 5 presents an empirical application to the cross-section of monthly U.S. stock returns. Section 6
summarizes the findings and concludes. The Appendix contains all the proofs.

## 2 The Model

We assume that there are a large number of securities, indexed by $i=1, \ldots, n$. Asset excess returns (returns minus the risk free rate) are observed for a number of time periods $t=1, \ldots, T$, where $n / T \rightarrow \infty$ as $n, T \rightarrow \infty$. We assume that the following characteristic-based factor model generates excess returns:

$$
\begin{equation*}
y_{i t}=f_{u t}+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}+\varepsilon_{i t}, \tag{1}
\end{equation*}
$$

where $y_{i t}$ is the excess return to security $i$ at time $t ; f_{u t}, f_{j t}$ are the factor returns; $g_{j}\left(X_{j i}\right)$ are the factor betas, where $X_{j i}$ are observable security characteristics. The factor returns $f_{j t}$ are linked to the security characteristics by the characteristic-beta functions $g_{j}(\cdot)$, which map characteristics to the associated factor betas. We assume that each $g_{j}(\cdot)$ is a smooth time-invariant function of characteristic $j$, but we do not assume a particular functional form. This generalizes the Rosenberg (1974) model where $g_{j}\left(X_{j i}\right)=X_{j i}$. The error terms in (1) are the asset-specific or idiosyncratic returns and they satisfy the conditional moment restriction that $E(\varepsilon \mid X, f)=0$, where $\varepsilon, X, f$ are vectors containing all the observations on the corresponding quantity. We shall suppose that all factors $f$ are unobserved, but our framework can easily accomodate the case where some factors are observed. This is essentially the model proposed by Connor and Linton (2007), although they only treated the case with fixed $T$. To simplify the exposition we are assuming that the characteristics $X_{j i}$ are time invariant. We will later on discuss the case where characteristics are allowed to vary over time.

The market factor $f_{u t}$ captures that part of common return not related to the security characteristics; all assets have unit beta to this factor. This factor captures the tendency of all equities to move together, irrespective of their characteristics. It is a common element in panel data models, see Hsiao (2003, section 3.6.2). In applications to returns data it is convenient to exclude own-effect intercept terms from (1) since they provide little benefit in terms of explanatory power and necessitate an additional time-series estimation step; see Connor and Korajczyk (1988, 1993), Connor and Linton (2007).

We also the allow for the case where one or more of the factors are directly observed and need not be estimated. So for example in our empirical analysis we investigate the special case where the market factor $f_{u t}$ is exogenously observable, and in particular is equal to a capitalization-weighted index. Other observable factors, such as the time-series innovations to unemployment or inflation,
or fixed income portfolio returns, are possible within our econometric framework.
Note that for fixed $t$, equation (1) constitutes a weighted additive nonparametric regression model for panel data, where the factor returns $f_{j t}$ are 'parametric weights' and the characteristic-beta functions $g_{j}(\cdot)$ are univariate nonparametric functions. Some discussion of additive nonparametric models can be found in Linton and Nielsen (1995). The situation here is somewhat nonstandard, since the same regression equation (1) holds each time period, with parametric weights varying each time period and the characteristic-beta functions time-invariant. We extend the weighted nonparametric regression methodology to account for this feature of time-varying weights in a pooled time-series, cross-sectional model.

Our model can be thought of as a special case of the usual statistical factor model

$$
\begin{equation*}
y_{i t}=\sum_{j=1}^{J} \beta_{i j} f_{j t}+\varepsilon_{i t}, \tag{2}
\end{equation*}
$$

where the unobserved factor loadings $\beta_{i j}$ are unrestricted, Ross (1976). Connor and Koracyzk (1993) developed the asymptotic principal component method for estimation of the factors in the case where the cross-section is large but the time series is fixed. Recent work of Bai and Ng (2002) and Bai $(2003,2004)$ have provided analysis for this method for the case where both $n$ and $T$ are large. Bai (2003) establishes pointwise asymptotic normality for estimates of the factors (at rate $\sqrt{n}$ ) and the loadings (at rate $\sqrt{T}$ ) under weak assumptions regarding cross-sectional and temporal dependence. ${ }^{1}$ The nesting of our model within (2) could be used for specification testing. Note however that in the case where the covariates in (1) are time varying, this nesting no longer holds.

### 2.1 Factor Scale Identification Conditions

In the case in which both characteristic-beta functions and factor returns are estimated from the data there is an obvious scale indeterminacy in the model. We assume that the observed characteristic $J$-vectors of the assets $X_{i}, i=1, \ldots, n$ are independent and identically distributed across $i$. We impose the identifying restrictions that for each factor the cross-sectional average beta equals zero and the cross-sectional variance of beta equals one, that is, $E^{*}\left[g_{j}\left(X_{j i}\right)\right]=0$ and $\operatorname{var}^{*}\left[g_{j}\left(X_{j i}\right)\right]=1$, where $E^{*}$ and var* denote moments with respect to some distribution. Note that this does not restrict the return model since the additive semiparametric model (1) is invariant to this rescaling. The choice of distribution to use in the normalization affects the interpretation of the factors. The condition $\operatorname{var}^{*}\left[g_{j}\left(X_{j i}\right)\right]=1$ sets the magnitude of factor return $j$; the conditions $E^{*}\left[g_{j}\left(X_{j i}\right)\right]=0$ affects the interpretation of the intercept. If we use the population distribution, then $E^{*}\left[g_{j}\left(X_{j i}\right)\right]=0$ means

[^1]that the intercept can be interpreted as the return to the average asset in the infinite population of assets; if we use a capitalization-weighted population distribution, then $E^{*}\left[g_{j}\left(X_{j i}\right)\right]=0$ means that the intercept can be interpreted as the return to the capitalization-weighted average asset. We will consider both of these in our empirical implementation; for simplicity our econometric theory focuses on the case in which we use the unweighted population distribution.

### 2.2 Estimating Additive Nonparametric Mispricing Functions

A central concern in the asset pricing literature is the determination of the expected returns on assets and their relationship to the risk exposures of the assets. Note that the information set of investors includes the characteristics $X$. Taking investors' expectation of excess returns $y_{i t}$ (conditional on the information set) using (1):

$$
\begin{equation*}
E\left[y_{i t}\right]=E\left[f_{u t}\right]+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) E\left[f_{j t}\right], \tag{3}
\end{equation*}
$$

which is the standard multi-factor asset pricing model: expected excess returns are linear in factor betas. Hence our model as developed so far imposes the standard multi-factor pricing condition on expected excess returns.

Our methodology provides a new asset pricing test against a general nonparametric pricing alternative. Fama and French (1993) create factor mimicking portfolios from size and value-sorted portfolios and then test for characteristic-related mispricing based on a finer grid of value and sizesorted portfolios. This two-stage procedure leaves open the question whether there is a hidden "identification condition" when using the same characteristics to create mimicking portfolios and to test for mispricing. We show that this is in fact the case.

Adapting the Fama-French mispricing test to our additive nonparametric framework generates an explicit identification condition. The characteristic-mispricing functions are only identified up to an orthogonality condition relative to the characteristic-beta functions. This is because the same characteristics are used to identify the factor risk premia and factor model mispricing.

We assume that there are mispricing inefficiencies given by smooth additive univariate nonparametric functions $\alpha_{j}\left(X_{j i}\right)$ using the same characteristics $X_{j i}$ as in the factor model. ${ }^{2}$ The return generating process becomes:

$$
\begin{equation*}
y_{i t}=f_{u t}+\sum_{j=1}^{J} \alpha_{j}\left(X_{j i}\right)+g_{j}\left(X_{j i}\right) f_{j t}+\varepsilon_{i t} . \tag{4}
\end{equation*}
$$

[^2]In order for the functions $\alpha_{j}($.$) to be identified we must impose:$

$$
\begin{gather*}
E\left[\alpha_{j}\left(X_{j i}\right)\right]=0  \tag{5}\\
E\left[\alpha_{j}\left(X_{j i}\right) g_{j}\left(X_{j i}\right)\right]=0 . \tag{6}
\end{gather*}
$$

The mean-zero condition (5) is standard in additive nonparametric models, so that the intercept $f_{u t}$ can be identified. The condition (6) is necessary in order for the risk premia of each factor return to be identified. To see why this is so, suppose that we relax the identification condition, then for any constant $a$ we can replace $\alpha_{j}\left(X_{j i}\right)$ with $\alpha_{j}^{*}\left(X_{j i}\right)=\alpha_{j}\left(X_{j i}\right)+a g_{j}\left(X_{j i}\right)$ and $f_{j t}^{*}=f_{j t}-a$ and the fit of the model is exactly the same. This indeterminacy is only eliminated by imposing (6). The intuition for the condition is clear: sample mean return which is in the linear span of the characteristic-beta function must be treated as "factor risk premia" rather than "mispricing."

### 2.3 An Important Special Case: An Observed Market Factor

As mentioned above it is possible to avoid estimation of one or more of the factor returns by replacing them with observed proxies. An important case is replacing the estimated intercept $f_{u t}$ with a market index return $f_{m t}$. Let returns obey a $J+1$ statistical factor model:

$$
\begin{equation*}
y_{i t}=\sum_{j=0}^{J} \beta_{i j} f_{j t}+\varepsilon_{i t} \tag{7}
\end{equation*}
$$

and let the market index return $f_{m t}$ be some linear combination of the $J+1$ factor returns. By a suitable factor rotation we can rewrite (7) as:

$$
y_{i t}=\beta_{0 i}^{*} f_{m t}+\sum_{j=1}^{J} \beta_{i j}^{*} f_{j t}^{*}+\varepsilon_{i t}
$$

where $\left\{f_{j t}^{*}, j=1, \ldots, J\right\}$ is the set of the other $J$ factors after the rotation. In order to implement our approach we impose $\beta_{0 i}^{*}=1$ and that $\sum_{j=1}^{J} \beta_{i j}^{*} f_{j t}^{*}$ can be written in the additive nonparametric form shown in (1). This produces the same factor model formulation as (1) but with $f_{u t}=f_{m t}$ observed rather than estimated. Note that the condition $\beta_{0 i}^{*}=1$ is well supported empirically (see Fama and French (1993) and Connor and Linton (2007)) and that this does not require that the market betas of all assets equal one. Rather it requires that the market betas of assets can be written as one plus a linear combination of their non-market factor exposures (see (9) in the next subsection).

In the case of an identified market factor (using a capitalization-weighted index) the identification conditions $E^{*}\left[g_{j}\left(X_{j i}\right)\right]=0$ must use the capitalization-weighted population distribution. This ensures that if we take capitalization-weighted averages of both sides of (1) we get $f_{m t}=f_{m t}$.

### 2.4 Embedding the CAPM as a Testable Restriction

Consider the case described in the previous subsection in which $f_{u t}=f_{m t}$ an observed market index. In this subsection we will describe how to restrict that version of the model to encompass the CAPM. Assume (admittedly, with a loss of generality, and against some of the empirical evidence) that the factor returns are i.i.d. through time. In this case we can easily embed the Capital Asset Pricing Model inside (3) and test the restriction. The CAPM imposes the pricing restriction:

$$
\begin{equation*}
E\left[y_{i t}\right]=\frac{\operatorname{cov}\left(y_{i t}, f_{m t}\right)}{\operatorname{var}\left(f_{m t}\right)} E\left[f_{m t}\right] . \tag{8}
\end{equation*}
$$

Using (1) to derive market covariances, and then dividing by $\operatorname{var}\left(f_{m t}\right)$ :

$$
\begin{equation*}
\frac{\operatorname{cov}\left(y_{i t}, f_{m t}\right)}{\operatorname{var}\left(f_{m t}\right)}=1+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) \frac{\operatorname{cov}\left(f_{j t}, f_{m t}\right)}{\operatorname{var}\left(f_{m t}\right)} \tag{9}
\end{equation*}
$$

Since the functions $g_{j}\left(X_{j i}\right)$ can take arbitrary values, the CAPM (8) will hold if and only if:

$$
\begin{equation*}
E\left[f_{j t}\right]=\frac{\operatorname{cov}\left(f_{j t}, f_{m t}\right)}{\operatorname{var}\left(f_{m t}\right)} E\left[f_{m t}\right] \tag{10}
\end{equation*}
$$

which is an easily testable condition. We present empirical evidence in Section Five below.

## 3 Estimation Strategy

For simplicity of exposition we focus on the case in which all the factors are estimated and there are no mispricing functions $\alpha_{j}\left(Z_{i j}\right)$ included in the factor model. Connor and Linton (2007) propose to estimate the period by period conditional expectation of $y_{i t}$ given the characteristics $X_{1 i}, \ldots, X_{J i}$ at a grid of points and then to estimate the factors and beta functions at the same grid of points using an iterative algorithm based on bilinear regression. This approach works well enough when the cross-section is very large and when $J$ is small, like two in their case. However, it is inefficient in general and works poorly in practice when $J$ is larger than two. For this reason we develop an alternative estimation strategy that makes efficient use of the restrictions embodied in (1). We use kernel methodology instead of series or splines because: (a) vastly more research has been published about them than these other methods; (b) we are able to derive the pointwise distribution theory for the estimated functions for this case only.

In order to describe the statistical properties of our estimators we make some assumptions about the data generating process. For notational convenience, we treat in detail the case of a fully balanced panel, where the set of assets and the characteristics of each asset do not vary through time. (In subsection 3.6 below we describe the modifications necessary for the case of an unbalanced panel.)

We assume that each $\varepsilon_{i t}$ is a martingale difference sequence with finite conditional and unconditional variance.

### 3.1 Population Characterization

To motivate our estimation methodology we first define the parameters of interest through a population least squares criterion. This is one way of defining the quantities $f, g$ consistent with (1). It has two advantages: (a) it only requires the conditional moment restriction $E(\varepsilon \mid X, f)=0$ for consistency; (b) it usually implies an efficient procedure under i.i.d. normal error terms. The solution to this population problem is characterized by first order conditions; to derive estimators we mimic this population first order condition by a sample equivalent. For clarity, we just treat the case where all the factors are unknown - if some factors are known then they do not need to be chosen in the optimization below.

Consider the population criterion

$$
\begin{equation*}
Q_{T}(f, g)=\frac{1}{T} \sum_{t=1}^{T} E\left[\left\{y_{i t}-f_{u t}-\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}\right\}^{2}\right] . \tag{11}
\end{equation*}
$$

In this criterion, the expectation is taken over the distribution of returns and characteristics, treating the factors as fixed parameters that are to be chosen - we are thinking of the factors as an exogenous stochastic process. Under some conditions there may exist a limiting (as $T \rightarrow \infty$ ) criterion function $Q(f, g)$ but we do not require this. We minimize $Q_{T}(f, g)$ with respect to the factors $f$ (which contains $f_{u t}, f_{j t}$ for all $\left.j, t\right)$ and the functions $g=\left(g_{1}, \ldots, g_{J}\right)$ subject to the identifying restrictions $E^{*}\left[g_{j}\left(X_{j i}\right)\right]=0$ and $\operatorname{var}^{*}\left[g_{j}\left(X_{j i}\right)\right]=1$.

This minimization problem can be characterized by a set of first order conditions for $f, g$. For expositional purposes we shall divide the problem in two: an equation characterizing $f$ given known $g$, and an equation characterizing $g$ given $f$.

### 3.1.1 Characterization of the Factor Returns

First we solve for the minimization of (11) over $f_{u t}, f_{j t}$ for all $j, t$ given $g($.$) is known. Note that if the$ population of assets is treated as fixed rather than random, then (11) simply amounts to a collection of unrelated cross-sectional regression problems, one per time period. In this case the solution to the minimization problem is obviously period-by-period least squares regression. We now show that this intuition extends to our environment with a random population of assets rather than a fixed cross-section.

Taking the first derivatives of (11) with respect to $f_{u t}, f_{j t}$ and setting to zero, the first order conditions are (for each $t=1, \ldots, T$ ):

$$
\begin{gather*}
E\left[\left\{y_{i t}-f_{u t}-\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}\right\}\right]=0  \tag{12}\\
E\left[\left\{y_{i t}-f_{u t}-\sum_{k=1}^{J} g_{k}\left(X_{k i}\right) f_{k t}\right\} g_{j}\left(X_{j i}\right)\right]=0, j=1, \ldots, J . \tag{13}
\end{gather*}
$$

These equations are linear in $f$ given $g$. This delivers a linear system of $J+1$ equations in $J+1$ unknowns for each time period $t$. Letting $f_{t}=\left[f_{u t}, f_{1 t}, \ldots, f_{J t}\right]^{\top}$, $y_{t}=\left[y_{1 t}, \ldots, y_{n t}\right]^{\top}$, and $G\left(X_{i}\right)=$ $\left[1, g_{1}\left(X_{1 i}\right), \ldots, g_{J}\left(X_{J i}\right)\right]^{\top}$ we have $E\left[G\left(X_{i}\right) G\left(X_{i}\right)^{\top}\right] f_{t}=E\left[G\left(X_{i}\right) y_{t}\right]$. It follows that

$$
f_{t}=E\left[G\left(X_{i}\right) G\left(X_{i}\right)^{\top}\right]^{-1} E\left[G\left(X_{i}\right) y_{t}\right] \equiv A^{-1} b_{t}
$$

provided $A=E\left[G\left(X_{i}\right) G\left(X_{i}\right)^{\top}\right]$ is non-singular, which we assume to be the case.

### 3.1.2 The Characteristic-beta Functions

Next we turn to the characterization of $g$ given $f$. Consider the point-wise derivative of (11) with respect to $g_{j}\left(X_{j}\right)$ conditional on a fixed value of $X_{j}$ :

$$
\lim _{\delta \rightarrow 0} \frac{1}{T} \sum_{t=1}^{T} E\left[\left\{y_{i t}-f_{u t}-\left\{g_{j}\left(X_{j i}\right)+\delta\right\} f_{j t}-\sum_{k \neq j} g_{k}\left(X_{k i}\right) f_{k t}\right\}^{2} \mid X_{j i}=x_{j}\right] / \delta
$$

Setting this derivative to zero gives a system of first-order conditions defining the criterion-minimizing functions $g_{j}(x)$ at value $x_{j}$ :

$$
\begin{align*}
\frac{1}{T} \sum_{t=1}^{T} f_{j t} E\left[y_{i t} \mid X_{j i}=x_{j}\right]= & \frac{1}{T} \sum_{t=1}^{T} f_{j t} f_{u t}+g_{j}\left(x_{j}\right) \frac{1}{T} \sum_{t=1}^{T} f_{j t}^{2}  \tag{14}\\
& +\frac{1}{T} \sum_{t=1}^{T} \sum_{k \neq j} f_{j t} f_{k t} E\left[g_{k}\left(X_{k i}\right) \mid X_{j i}=x_{j}\right]
\end{align*}
$$

for $j=1, \ldots, J$. These equations are linear in $g$ given $f$ but they are only implicit, that is, they constitute a system of integral equations (of type 2) in the functional parameter $g$, see Mammen, Linton, and Nielsen (1999) and Linton and Mammen (2005). In the appendix we give some discussion of the properties of this integral equation. To simplify the notation define the conditional expectations:

$$
\lambda_{1 t}(j, x)=E\left[y_{i t} \mid X_{j i}=x\right] \quad ; \quad \lambda_{2}(j, k, x)=E\left[g_{k}\left(X_{k i}\right) \mid X_{j i}=x\right] .
$$

Substituting these functions into (14) and rearranging we obtain the system of equations

$$
\begin{equation*}
g_{j}\left(x_{j}\right)=\frac{\sum_{t=1}^{T} f_{j t}\left[\lambda_{1 t}\left(j, x_{j}\right)-f_{u t}-\sum_{k \neq j} f_{k t} \lambda_{2}\left(j, k, x_{j}\right)\right]}{\sum_{t=1}^{T} f_{j t}^{2}}, j=1, \ldots, J . \tag{15}
\end{equation*}
$$

The solution to (15) is typically not unique due to the lack of identification. One can impose the identifying restrictions by instead considering the constrained optimization problem and manipulating the first order condition of the associated Lagrangean. An equivalent approach is to take any unrestricted solution $g_{j}\left(x_{j}\right)$ and replace it by

$$
\begin{equation*}
\bar{g}_{j}\left(x_{j}\right)=\frac{g_{j}\left(x_{j}\right)-\int g_{j}\left(x_{j}\right) d P_{j}^{*}\left(x_{j}\right)}{\sqrt{\int g_{j}^{2}\left(x_{j}\right) d P_{j}^{*}\left(x_{j}\right)}} \tag{16}
\end{equation*}
$$

which also satisfies (15), where $P_{j}^{*}$ is the probability distribution associated with the identifying assumption on $g_{j}($.$) .$

### 3.2 Estimation

### 3.2.1 The Kernel Estimates

The first set of kernel estimates measure the period-specific expected return of an asset given its $j^{\text {th }}$ characteristic is $x_{j}$ and conditional on the observed factor returns in period $t$. We use the following (boundary corrected) kernel estimate:

$$
\widehat{\lambda}_{1 t}\left(j, x_{j}\right)=\frac{\sum_{i=1}^{n} K_{h}\left(X_{j i}, x_{j}\right) y_{i t}}{\sum_{i=1}^{n} K_{h}\left(X_{j i}, x_{j}\right)}
$$

where for each $x$ in the support of $X_{t}, K_{h}(x, y)=K_{h}^{x}(x-y)$ for some kernel $K^{x}$ such that $K_{h}^{x}(u)=$ $h^{-1} K^{x}\left(h^{-1} u\right)$ and $K_{h}^{x}(u)=K_{h}(u)$ for all $x$ in the interior of the support of $X_{j i}$. Here, $K_{h}()=$. $K(. / h) / h$ and $K$ is a kernel while $h$ is a bandwidth. We shall assume that each covariate is compactly supported, which is why we need this boundary adjustment.

The second set of kernel estimates give the expected factor beta of asset $i$ for factor $j$ based on the asset's characteristic $X_{j i}$ and the estimated characteristic-beta function $g_{k}^{[i]}(\cdot)$ :

$$
\widehat{\lambda}_{2}^{[i]}\left(j, k, x_{j}\right)=\frac{\sum_{i=1}^{n} K_{h}\left(X_{j i}, x_{j}\right) g_{k}^{[i]}\left(X_{k i}\right)}{\sum_{i=1}^{n} K_{h}\left(x_{j}, x_{j}\right)}
$$

Both of these kernel estimates are familiar features of weighted additive nonparametric regression. Note that $\hat{\lambda}_{2}^{[i]}$ does not require a time subscript since under our assumption of a fully-balanced panel, all assets have constant characteristics over time. We will weaken this assumption later.

### 3.2.2 Estimation of Factor Returns and Characteristic-Beta Functions

We replace the unknown quantities in $A, b_{t}$ and equation (15) by estimated values, denoted by hats, and iterate between the factor return $f$ and characteristic-beta function $g($.$) estimation problems.$ The solution for $f$ depends upon $g($.$) , and the solution for g_{j}($.$) depends both upon f$ and $\left\{g_{k}(\right.$.$) ,$ $k \neq j\}$. We use the Gauss-Seidel iteration to reconcile these component solutions.

1. Let $f^{[0]}, g^{[0]}($.$) be initial estimates.$
2. Then let for all $x$ and $i=1,2, \ldots$, let

$$
\begin{align*}
\widehat{g}_{j}^{[i+1]}(x)= & \frac{\sum_{t=1}^{T} \widehat{f}_{j t}^{[i]}\left(\widehat{\lambda}_{1 t}(j, x)-\widehat{f}_{u t}^{[i]}\right)}{\sum_{t=1}^{T} \widehat{f}_{j t}^{[i] 2}}-\frac{\sum_{t=1}^{T} \sum_{k>j} \widehat{f}_{j t}^{[i]} \widehat{f}_{k t}^{[i]} \hat{\lambda}_{2}^{[i]}(j, k, x)}{\sum_{t=1}^{T} \widehat{f}_{j t}^{[i]}}  \tag{17}\\
& -\frac{\sum_{t=1}^{T} \sum_{k<j} \widehat{f}_{j t}^{[i]} \hat{f}_{k t}^{[i]} \hat{\lambda}_{2}^{[i+1]}(j, k, x)}{\sum_{t=1}^{T} \hat{f}_{j t}^{[i] 2}} \\
\widehat{g}_{j}^{[i+1]}(x)= & \frac{\widehat{g}_{j}^{[i+1]}(x)-\int \widehat{g}_{j}^{[i+1]}(x) d P_{j}^{*}(x)}{\sqrt{\int \widehat{g}_{j}^{[i+1]}(x)^{2} d P_{j}^{*}(x)}}, \tag{18}
\end{align*}
$$

where hats denote estimated quantities defined above.
3. Then given estimates $\widehat{g}_{j}^{[i]}\left(X_{j i}\right)$ from the previous iteration on $g$ given $f$ we compute for each $t$

$$
\begin{equation*}
\widehat{f}_{t}^{[i+1]}=\left[\frac{1}{n} \sum_{i=1}^{n} \widehat{G}^{[i]}\left(X_{i}\right) \widehat{G}^{[i]}\left(X_{i}\right)^{\top}\right]^{-1} \frac{1}{n} \sum_{i=1}^{n} \widehat{G}^{[i]}\left(X_{i}\right) y_{i t}, \tag{19}
\end{equation*}
$$

where $\widehat{G}^{[i]}\left(X_{i}\right)=\left[1, \widehat{g}_{1}^{[i]}\left(X_{1 i}\right), \ldots, \widehat{g}_{J}^{[i]}\left(X_{J i}\right)\right]^{\top}$.
One can iterate (17),(18), and (19) repeatedly until some convergence criterion is met. The convergence properties of this algorithm are not studied here, we refer the reader to Mammen, Linton, and Nielsen (1999) for a fuller discussion of this issue in a special case of our model.

### 3.2.3 Initial Estimators

We describe two approaches to finding starting values. The first approach produces consistent starting values under the model assumptions, and is based on using time averaged data. This approach has similarities to the averaging method proposed in Pesaran (2006) except that our averaging is over time rather than over the cross-section. ${ }^{3}$ In particular, let $\left\{w_{T t}\right\}$ be some sequence with

[^3]$\sum_{t=1}^{T} w_{T t} \leq \bar{w}<\infty$, and let $\bar{y}_{i}=\sum_{t=1}^{T} w_{T t} y_{i t}, \bar{\varepsilon}_{i}=\sum_{t=1}^{T} w_{T t} \varepsilon_{i t}, \bar{f}_{u}=\sum_{t=1}^{T} w_{T t} f_{u t}, \bar{f}_{j}=\sum_{t=1}^{T} w_{T t} f_{j t}$. For example, $w_{T t}=1 / T$. Then we have
\[

$$
\begin{equation*}
\bar{y}_{i}=\bar{f}_{u}+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) \bar{f}_{j}+\bar{\varepsilon}_{i}=\bar{f}_{u}+\sum_{j=1}^{J} \bar{g}_{j}\left(X_{j i}\right)+\bar{\varepsilon}_{i}, \tag{20}
\end{equation*}
$$

\]

where $\bar{g}_{j}()=.g_{j}(.) \bar{f}_{j}$. This constitutes an additive nonparametric regression with components $\bar{g}_{j}($. that are mean zero, i.e., $E\left[\bar{\varepsilon}_{i} \mid X_{i}, f\right]=0$ and $E\left[\bar{g}_{j}\left(X_{j i}\right)\right]=0, j=1, \ldots, J$. Then, since $g_{j}($.$) is propor-$ tional to $\bar{g}_{j}($.$) we see that in fact g_{j}()=.\bar{g}_{j}(.) / \int \bar{g}_{j}\left(x_{j}\right)^{2} d P_{j}^{*}\left(x_{j}\right)$, i.e., a corresponding renormalization of $\bar{g}_{j}$ yields $g_{j}$. This means that we can estimate the functions $\bar{g}_{j}($.$) by the smooth backfitting method$ of Mammen, Linton, and Nielsen (1999) and then apply an empirical renormalization to obtain estimates of $g_{j}$. The quantity $\bar{f}_{u}$ can be estimated by the grand mean $\bar{y}=\sum_{i=1}^{n} \bar{y}_{i} / n$. The theory of Mammen, Linton, and Nielsen (1999) can be directly applied to the estimator of $\bar{g}_{j}$ provided that $\bar{f}_{j} \neq 0$, with some slight differences. Formally, the setting is that of a triangular array because the distribution of $\bar{\varepsilon}_{i}$ changes with $T(n)$, and in fact the error term in (20) is $O_{p}\left(T^{-1 / 2}\right) .{ }^{4}$ This makes the convergence rate of the nonparametric estimators faster by this magnitude. To estimate the factors $f_{t}$ rather than their time series means we must go back and cross-sectionally regress $y_{i t}$ on a constant and $\widehat{g}_{1}\left(X_{1 i}\right), \ldots, \widehat{g}_{J}\left(X_{J i}\right)$ for each time period $t .{ }^{5}$ The procedure described above provides consistent initial estimates for the case of the fully balanced panel.

The second approach to starting values uses a variant of Rosenberg's (1974) linear model:

$$
\begin{equation*}
g_{j}\left(X_{j i}\right)=X_{j i} . \tag{21}
\end{equation*}
$$

In this linear case it is simple to rescale the characteristics so that the identification constraints hold using (21). We scale the mean and variance of the characteristics so that $E^{*}\left[X_{j i}\right]=0$ and $\operatorname{var}^{*}\left[X_{j i}\right]=1$ for each $j$; for each characteristic, this just requires subtracting the weighted crosssectional mean and dividing by the weighted cross-sectional standard deviation each time period. The simple linear model for $g($.$) gives rise to a linear cross-sectional regression model to estimate$ $f_{j t}$ :

$$
\begin{equation*}
y_{i t}=f_{u t}+\sum_{j=1}^{J} X_{j i} f_{j t}+\varepsilon_{i t} \tag{22}
\end{equation*}
$$

We begin with ordinary least squares estimation of (22). These estimates of $f_{u t}$ and $f_{j t}$ serve merely as starting values and have no consistency properties. Connor and Linton (2007) find that this linear model provides quite a reasonable first approximation. As long as these initial estimates are in a convergent neighborhood of the maximizing values, the final estimates after repeated iteration will be unaffected. We use this second approach in our empirical implementation below.

[^4]
### 3.3 Unbalanced, Time-varying Panel Data

We have so far mostly assumed a fully balanced panel dataset. The set of observed assets is assumed constant over time, with each asset having a fixed vector of characteristic betas. The only time variation in this fully balanced panel comes through the random factor realizations and random asset-specific returns. In applications, the set of assets must be allowed to vary over the time sample, since the set of equities with full records over a reasonably long sample period is a small subset of the full dataset. Also, the characteristics of the assets must be allowed to vary through time.

We assume that the observations are unbalanced in the sense that in time period $t$ we only observe $n_{t}$ firms (for simplicity labelled $i=1, \ldots, n_{t}$ ), but that the source of "missingness" is random. Also, we assume that the characteristics are time varying but stationary over time for each $i$ and i.i.d. over $i$. This yields first order conditions for $f, g$ that are similar to the balanced case. Now the matrix $E\left[G\left(X_{i t}\right) G\left(X_{i t}\right)^{\top}\right]$ for example depends on time, while the expression for $g_{j}$ becomes

$$
\begin{equation*}
g_{j}\left(x_{j}\right)=\frac{\frac{1}{T} \sum_{t=1}^{T} \frac{1}{n_{t}} \sum_{i=1}^{n_{t}} f_{j t}\left(E\left[y_{i t} \mid X_{j i t}=x_{j}\right]-f_{u t}-\sum_{k \neq j} f_{k t} E\left[g_{k}\left(X_{k i t}\right) \mid X_{j i t}=x_{j}\right]\right)}{\frac{1}{T} \sum_{t=1}^{T} f_{j t}^{2}} \tag{23}
\end{equation*}
$$

The estimation algorithm is essentially the same as outlined above. The two unknown expectations in (23) are replaced by kernel estimates given by

$$
\begin{align*}
\widehat{\lambda}_{1 t}\left(j, x_{j}\right) & =\frac{\sum_{i=1}^{n_{t}} K_{h}\left(X_{j i t}, x_{j}\right) y_{i t}}{\sum_{i=1}^{n_{t}} K_{h}\left(X_{j i t}, x_{j}\right)} \\
\widehat{\lambda}_{2 t}^{[i]}\left(j, k, x_{j}\right) & =\frac{\sum_{i=1}^{n_{t}} K_{h}\left(X_{j i t}, x_{j}\right) g_{k}^{[i]}\left(X_{k i t}\right)}{\sum_{i=1}^{n_{t}} K_{h}\left(X_{j i t}, x_{j}\right)} . \tag{24}
\end{align*}
$$

Note that $\widehat{\lambda}_{1 t}$ and $\widehat{\lambda}_{2 t}^{[i]}\left(j, k, x_{j}\right)$ now both have to be estimated in each sample period separately to allow for time variation.

## 4 Distribution Theory

In this section we provide the distribution theory for our estimates of the factors and of the characteristic functions in the balanced case. The theorem covers the estimators defined by a finite number of iterations of (17),(18), and (19) starting from the initial consistent estimators described under (20). The general approach for obtaining the asymptotic distribution uses the methods developed in Mammen, Linton, and Nielsen (1999) and Linton and Mammen (2005) for treating estimators defined as the solutions of type 2 linear integral equations. The novelty here is due to the weighting by the factors and the fact that we wish to allow both the cross-section and the time dimension
to grow. Regarding the asymptotics, we take joint limits as $n, T \rightarrow \infty$ under the restriction that $n / T \rightarrow \infty$ as described in Phillips and Moon (1999, Definition 2(b)).

Let $\mathcal{F}_{a}^{b}$ be the $\sigma$-algebra of events generated by the vector random variable $\left\{U_{t} ; \quad a \leq t \leq b\right\}$. The processes $\left\{U_{t}\right\}$ is called strongly mixing [Rosenblatt (1956)] if

$$
\sup _{1 \leq t} \sup _{A \in \mathcal{F}_{-\infty}^{t}, B \in \mathcal{F}_{t+k}^{\infty}}|\operatorname{Pr}(A \cap B)-\operatorname{Pr}(A) \operatorname{Pr}(B)| \equiv \alpha(k) \rightarrow 0 \quad \text { as } k \rightarrow \infty .
$$

We make the following assumptions.
Assumptions A.
A1. The double array $\left\{X_{i}, \varepsilon_{i t}\right\}_{i, t=1}^{n, T}$ are defined on a probability space $(\Omega, \mathcal{F}, P)$. For each $t,\left(X_{i}, \varepsilon_{i t}\right)$ are i.i.d. across $i$. The processes $\left\{\varepsilon_{i t}\right\}$ are strongly mixing with mixing coefficients, $\alpha_{i}(k)$ such that $\alpha_{i}(k) \leq C \bar{\alpha}^{k}$ for some $C \geq 0$ and some $\bar{\alpha}<1$. Furthermore, $E\left(\varepsilon_{i t} \mid \mathcal{F}_{t-1}\right)=0$ and $E\left(\varepsilon_{i t}^{2} \mid \mathcal{F}_{t-1}\right)=\sigma_{t}^{2}\left(X_{i}\right)$ a.s., where $\mathcal{F}_{t-1}$ is the sigma field generated by $X_{i}$ and the past of $\varepsilon_{i t}$. Furthermore, for some $\kappa>4$, $\sup _{t} E\left[\left|\varepsilon_{i t}\right|^{\kappa}\right]<\infty$.

A2. The covariate $X_{i}=\left(X_{1 i}, \ldots, X_{J i}\right)^{\top}$ has absolutely continuous density $p$ supported on $\mathcal{X}=$ $[\underline{x}, \bar{x}]^{J}$ for some $-\infty<\underline{x}<\bar{x}<\infty$. The functions $g_{j}(\cdot)$ together with the density $p(\cdot)$ are twice continuously differentiable over the interior of $\mathcal{X}$ and are bounded on $\mathcal{X}$. The density function $p(x)$ is strictly positive at each $x \in \mathcal{X}$. Denote by $p_{j}(x)$ the marginal probability density for characteristic $j$ with support $\mathcal{X}_{j}=[\underline{x}, \bar{x}]$. The matrix $E\left[G\left(X_{i}\right) G\left(X_{i}\right)^{\top}\right]$ is strictly positive definite.

A3. For each $x \in[\underline{x}, \bar{x}]$ the kernel function $K^{x}$ has support $[-1,1]$ and satisfies $\int K^{x}(u) d u=1$ and $\int K^{x}(u) u d u=0$, such that for some constant $C, \sup _{x \in[\underline{x}, \bar{x}]}\left|K^{x}(u)-K^{x}(v)\right| \leq C|u-v|$ for all $u, v \in[-1,1]$. Define $\mu_{j}(K)=\int u^{j} K(u) d u$ and $\|K\|_{2}^{2}=\int K^{2}(u) d u$. The kernel $K$ is bounded, has compact support ( $\left[-c_{1}, c_{1}\right]$, say), is symmetric about zero, and is Lipschitz continuous, i.e., there exists a positive finite constant $C_{2}$ such that $|K(u)-K(v)| \leq C_{2}|u-v|$.

A4. $n, T \rightarrow \infty$ in such a way that $n / T \rightarrow \infty$.
A5. For some $a \geq 2$, the quantities $\sup _{T \geq 1} \sum_{t=1}^{T} f_{j t}^{a} / T<\infty$. The quantities $\Phi_{j}=\lim _{T \rightarrow \infty} T^{-1} \sum_{t=1}^{T} f_{j t}^{2}$ and $\Psi_{j}\left(x_{j}\right)=\lim _{T \rightarrow \infty} T^{-1} \sum_{t=1}^{T} f_{j t}^{2} \sigma_{j t}^{2}\left(x_{j}\right)$ exist and $\Phi_{j}>0$. Here, $\sigma_{j t}^{2}\left(x_{j}\right)=E\left[\varepsilon_{i t}^{2} \mid X_{j i}=x_{j}\right]$.

A6. The bandwidth sequence $h(n, T)$ satisfies $n h^{4}(\log T)^{\rho} \rightarrow 0$ and $n T h \rightarrow \infty$ as $n, T \rightarrow \infty$.
A7. For $j=u, 1, \ldots, J$, there exists $\rho^{\prime}, C \in(0, \infty)$ such that $\max _{1 \leq t \leq T}\left|f_{j t}\right| \leq C(\log T)^{\rho^{\prime}}$.

In A1 we allow $\varepsilon_{i t}$ to have certain types of nonstationarity: this is possible since the CLT is coming from the cross-sectional independence. In particular, we can allow general form of time series and cross-sectional conditional heteroskedasticity. We could also allow $\varepsilon_{i t}$ to be weakly autocorrelated, with no change in the limiting distributions, but the data do not seem to require this, so for simplicity we maintain mean unpredictability of the errors. We can also allow $\varepsilon_{i t}$ to be cross-sectionally correlated without affecting the result for $\widehat{g}_{j}$. However, this will affect the asymptotic distribution of $\widehat{f}_{t}$ in such a way that it makes inference difficult without some assumptions on the nature of this cross-sectional dependence. Bai (2003) allows for cross-sectional dependence in his main results but his inferential results for the estimated factors assume cross-sectional independence for the same reason. In A2 we only require second order smoothness of the functions $g_{j}$ regardless of how many of them there are. In this sense our method completely avoids the curse of dimensionality, Horowitz and Mammen (2006). Assumption A5 embodies an assumption about the magnitudes of the factors; here we assume that they behave like the outcome of a stationary process with finite moments of order $a$. This could be relaxed to allow either faster growth in $\sum_{t=1}^{T} f_{j t}^{2}$, reflecting nonstationary factors, or slower growth reflecting perhaps many zero values in the factors, but we do not do this here as the data seem to support our assumption. Assumption A7 is needed for the uniform convergence rates below. This assumption is consistent with the factors being realizations from a stationary Gaussian process. Again, this condition could be weakened to allow faster growth in $\max _{1 \leq t \leq T}\left|f_{j t}\right|$ at the expense of further restrictions elsewhere.

Define:

$$
\begin{gather*}
\Omega_{j}\left(x_{j}\right)=\frac{\Psi_{j}\left(x_{j}\right)}{p_{j}\left(x_{j}\right) \Phi_{j}^{2}}\|K\|_{2}^{2}  \tag{25}\\
V_{t, t}=E\left[G\left(X_{i}\right) G\left(X_{i}\right)^{\top}\right]^{-1} E\left[\varepsilon_{i t}^{2} G\left(X_{i}\right) G\left(X_{i}\right)^{\top}\right] E\left[G\left(X_{i}\right) G\left(X_{i}\right)^{\top}\right]^{-1} . \tag{26}
\end{gather*}
$$

Theorem 1. Suppose that Assumptions A1-A6 hold. Then, for each $t$ there exists a stochastically bounded sequence $\delta_{n, t}$ such that

$$
\begin{equation*}
\sqrt{n}\left(\widehat{f_{t}}-f_{t}-h^{2} \delta_{n, t}\right) \Longrightarrow N\left(0, V_{t, t}\right) \tag{27}
\end{equation*}
$$

while $\widehat{f}_{t}, \widehat{f}_{s}$ with $t \neq s$ are asymptotically independent. Suppose also that $A 7$ holds. Then, for some $\rho>0$

$$
\begin{equation*}
\max _{1 \leq t \leq T}\left|\widehat{f}_{t}-f_{t}\right|=O_{p}\left(\left(n^{-1 / 2}+h^{2}\right)(\log T)^{\rho}\right) \tag{28}
\end{equation*}
$$

Also, given Assumptions A1-A7, there exists a bounded continuous function $\beta_{j}($.$) such that for$ each $x_{j} \in(\underline{x}, \bar{x})$,

$$
\begin{equation*}
\sqrt{n T h}\left(\widehat{g}_{j}\left(x_{j}\right)-g_{j}\left(x_{j}\right)-h^{2} \beta_{j}\left(x_{j}\right)\right) \Longrightarrow N\left(0, \Omega_{j}\left(x_{j}\right)\right) . \tag{29}
\end{equation*}
$$

Furthermore, $\widehat{g}_{j}\left(x_{j}\right), \widehat{g}_{k}\left(x_{k}\right)$ are asymptotically independent for any $x_{j}, x_{k}$ for $j \neq k$.

## Remarks.

1. Provided the bandwidth is chosen so that $n h^{4} \rightarrow 0$, the estimator $\widehat{f_{t}}$ is consistent at rate $n^{-1 / 2}$ and the asymptotic distribution is as if the characteristic functions were known and least squares were applied. The estimators $\widehat{g}_{j}\left(x_{j}\right)$ are consistent at rate $(n T)^{-2 / 5}$ provided a bandwidth of order $(n T)^{-1 / 5}$ is chosen and under some restrictions on the rates at which $T, n$ increase. This should be the optimal rate for this problem, Stone (1980). ${ }^{6}$ It can be that $\widehat{g}_{j}\left(x_{j}\right)$ converges to $g_{j}\left(x_{j}\right)$ faster than $\widehat{f_{t}}$ converges to $f_{t}$; this happens when $T^{4} / n \rightarrow \infty$. This is because of the extra pooling over time in the specification of $g_{j}$. Note that the asymptotic variance of the characteristic function estimates is as if the factors were known. Finally, the estimated factors are asymptotically independent of the estimated functions.
2. Suppose that $\varepsilon_{i t}$ is homoskedastic, i.e., $\sigma_{j t}^{2}\left(x_{j}\right)=\sigma_{\varepsilon}^{2}$ for all $j$, $t$, then the asymptotic variance of $\widehat{g}_{j}\left(x_{j}\right)$ simplifies to $\Omega_{j}\left(x_{j}\right)=\left(\sigma_{\varepsilon}^{2} /\left(p_{j}\left(x_{j}\right) \Phi_{j}\right)\right)\|K\|_{2}^{2}$. We argue that this is a natural 'oracle' bound for the performance of estimators in this case, see Linton (1997). Suppose that we could observe the partial residuals $U_{j i t}=y_{i t}-f_{u t}-\sum_{k \neq j} f_{k t} g_{k}\left(X_{k i}\right)$ and the factors, then we can compute the pooled regression smoother

$$
\begin{equation*}
\widehat{g}_{j}^{\text {oracle }}(x)=\frac{\sum_{t=1}^{T} \sum_{i=1}^{n} K_{h}\left(X_{j i}, x_{j}\right) f_{j t} U_{j i t}}{\sum_{t=1}^{T} \sum_{i=1}^{n} K_{h}\left(X_{j i}, x_{j}\right) f_{j t}^{2}} . \tag{30}
\end{equation*}
$$

This shares the asymptotic variance of our estimator, and, since it uses more information than we have available, it is comforting that our estimator performs as well as it. We note that it should also be possible to establish, along the lines of Horowitz, Klemela, and Mammen (2006), that a local linear implementation of $\widehat{g}_{j}^{\text {oracle }}(x)$ is Best Linear Minimax.
3. Standard errors can be obtained in an obvious way by plugging in estimated quantities. In particular, valid standard errors for the factors can be obtained from the final stage least squares regression of returns on the characteristic functions. Standard errors for $\widehat{g}_{j}\left(x_{j}\right)$ can be computed from

$$
\begin{equation*}
\frac{\widehat{\Omega_{j}\left(x_{j}\right)}}{n T h}=\frac{\sum_{t=1}^{T} \sum_{i=1}^{n} K_{h}^{2}\left(X_{j i}, x_{j}\right) \widehat{f}_{j t}^{2} \widehat{\epsilon}_{i t}^{2}}{\left(\sum_{t=1}^{T} \sum_{i=1}^{n} K_{h}\left(X_{j i}, x_{j}\right) \widehat{f}_{j t}^{2}\right)^{2}}, \tag{31}
\end{equation*}
$$

where $\widehat{\varepsilon}_{i t}=y_{i t}-\widehat{f}_{u t}-\sum_{j=1}^{J} \widehat{f}_{j t} \widehat{g}_{j}\left(X_{j i}\right)$ are residuals computed from the estimated factors and characteristic functions, see Fan and Yao (1998) for discussion of nonparametric standard errors.
4. Bandwidth selection can be handled in the framework of penalized least squares, see Mammen

[^5]and Park (2005). Adapted to our problem, this involves choosing bandwidth $h$ to minimize
$$
P L S(h)=\frac{1}{n T} \sum_{t=1}^{T} \sum_{i=1}^{n} \widehat{\varepsilon}_{i t}^{2}\left[1+\frac{2 J K(0)}{n T h}\right]
$$
over a set $H$ of bandwidths. Mammen and Park showed the consistency of this method for additive nonparametric regression in a single cross-section. In practice, we have found it important to allow the bandwidth to vary with location, and we describe below an alternative more pragmatic approach to bandwidth selection.
5. In practice it might be important to determine the number of factors to include. Bai and Ng (2002) propose a test for the inclusion of an additional factor based on residual sum of squares, which can be adapted to our framework. Specifically, let
$$
P C(J)=\frac{1}{n T} \sum_{t=1}^{T} \sum_{i=1}^{n} \widehat{\varepsilon}_{i t}^{2}(J)+J \varkappa(n, T),
$$
where $\varkappa(n, T) \rightarrow 0$ and $\min \{n, T\} \varkappa(n, T) \rightarrow \infty$. Then we expect that $\operatorname{Pr}[\widehat{J}=J] \rightarrow 1$ as $n, T \rightarrow \infty$, where $\widehat{J}=\arg \min _{J \leq J_{\max }} P C(J)$ and $J_{\max }$ is some predetermined upper bound.
6. The results (27)-(29) follow also for the unbalanced case with suitable generalizations. The case where the covariate process is stationary is particularly simple because then one only needs to replace $n$ by $n_{t}$ in (27), $n$ by $\min _{1 \leq t \leq T} n_{t}$ in (28), and $n T$ by $\sum_{t=1}^{T} n_{t}$ in (29).
7. Our model imposes two main testable restrictions on the conditional expectation $E\left[y_{i t} \mid X_{i}\right]=$ $m_{t}\left(X_{i}\right)$ : poolability, and additivity. ${ }^{7}$ Additivity refers to the separability of the covariate effect. This structure is present in nearly all regression models for stock returns: it is essential for the conventional interpretation of risk factors. Poolability refers to the time invariance of the functions $g_{j}($.$) . This$ restriction is also very convenient for interpretation and forecasting. Baltagi, Hidalgo, and Li (1996) propose a test of poolability that can be adapted to our framework based on comparing pooled and unpooled estimators. Gozalo and Linton (2001) have proposed tests of additivity in a crosssectional setting based on comparing restricted with unrestricted estimators that work with marginal integration estimators, Linton and Nielsen (1995). ${ }^{8}$ We do not pursue this in our application because applied researchers are more likely to question why we need all the complicated nonparametric machinery, so we next turn to testing a given parametric shape against our more general model. This can be done formally using the test statistics developed for example by Härdle and Mammen (1993), see Li and Racine (2007). We have informally investigated this issue in our application using

[^6]the pointwise standard errors, but since these are tiny, and the shapes are nonlinear, the formal test would clearly reject.

### 4.1 Time Series Analysis

In order to test the embedded CAPM pricing restriction (10) and to examine the dynamic behaviour of the factors, we need to apply the above theory to modelling the time series behaviour of the estimated factors. A large literature has considered this problem for specific models, including Stock and Watson (1998) and Bai (2003). Suppose that the true factors obey some time series model that is consistent with our earlier treatment of the factors as fixed under the assumption of strong exogeneity. In particular, suppose that

$$
\begin{equation*}
E\left[\psi\left(F_{t}, Z_{t} ; \theta_{0}\right)\right]=0 \tag{32}
\end{equation*}
$$

for some observed instruments $Z_{t}$ for some true value $\theta_{0}$ of a vector of parameters $\theta \in \mathbb{R}^{p}$, where $\psi$ is a $q$-vector of moment conditions with $q \geq p$ and $F_{t}=\operatorname{vec}\left(f_{t}, f_{t-1}, \ldots, f_{t-k}\right)$. This includes a number of cases of interest. In our empirical application we consider the case in which $f_{t}$ follows a vector autoregression $A(L) f_{t}=\eta_{t}$, where $\eta_{t}$ is i.i.d. and $A(L)=A_{0}-A_{1} L-\cdots-A_{r} L^{r}$ for parameter matrices $A_{0}, A_{1}, \ldots, A_{r}$, see Borak, Härdle, Mammen, and Park (2007). It follows that $E\left[A(L) f_{t} \otimes Z_{t}\right]=0$ for any $Z_{t}$ in the past of $f_{t}$.

To estimate the parameters we use the estimated factors and minimize the quadratic form

$$
\begin{gather*}
\widehat{M}_{T}(\theta)^{\top} W_{T} \widehat{M}_{T}(\theta)  \tag{33}\\
\widehat{M}_{T}(\theta)=\frac{1}{T-k} \sum_{t=k+1}^{T} \psi\left(\widehat{F}_{t}, Z_{t} ; \theta\right)
\end{gather*}
$$

with respect to $\theta$, where $W_{T}$ is a symmetric positive definite weighting matrix. Let $\widehat{\theta}$ be any minimizer of (33). For simplicity we assume that the factors are stationary and mixing; for most finance applications the assumptions of stationarity (or at least local stationarity) and mixing seem reasonable.

Hansen, Nielsen, and Nielsen (2004) consider the problem of using estimated values $\widehat{f_{t}}$ in linear time series models for an unobserved quantity $f_{t}$. They prove a general result that provided $\sum_{t=1}^{T}\left(\widehat{f_{t}}-\right.$ $\left.f_{t}\right)^{2} \xrightarrow{P} 0$ as $T \rightarrow \infty$, then we may use the predicted time series as if it was the true unobserved time series for instance in estimation and unit root testing in the sense that using the estimated values leads to the same asymptotic distribution (for $T \rightarrow \infty$ ) as if the true values were used. Their result applies to linear models for stationary and nonstationary factors. However, in the case of stationary factors where $\sqrt{T}$ consistent estimation of $\theta_{0}$ is possible, this condition is too strong. Specifically, it suffices
that there is a linear expansion for $\widehat{f_{t}}-f_{t}$ and that the uniform rate $\max _{1 \leq t \leq T}\left|\widehat{f}_{t}-f_{t}\right|=o_{p}\left(T^{-1 / 4}\right)$ holds, both of which are obtainable from Theorem 1. In addition, our setting allows for quite general nonlinearity in the time series model.

For the factor modelling result we treat the factors as random and we need some additional assumptions. We suppose that $\left\{f_{t}, Z_{t}\right\}_{t=1}^{\infty}$ is a jointly stationary process satisfying strong mixing. In this case one should interpret $\left\{f_{t}, Z_{t}\right\}_{t=1}^{\infty}$ as being independent of $X, \varepsilon$ and A5 and A7 as holding with probability one. Assumption A7 holds with probability one when $f_{t}$ is a stationary mixing Gaussian process.

## Assumptions B.

B1. We suppose that the process $\left\{\left(f_{t}, Z_{t}\right)\right\}$ is strictly stationary and strong mixing with $\alpha(k)$ satisfying $\alpha(k) \leq C \bar{\alpha}^{k}$ for some $C \geq 0$ and some $\bar{\alpha}<1$.

B2. The limiting moment condition $M(\theta)=E\left[\psi\left(F_{t}, Z_{t} ; \theta\right)\right]$ has a unique zero at $\theta=\theta_{0}$, where $\theta_{0}$ is an interior point of the compact parameter set $\Theta \subseteq \mathbb{R}^{p}$. Furthermore, $\psi\left(F_{t}, Z_{t} ; \theta_{0}\right)$ is a martingale difference sequence.

B3. The function $\psi$ is twice continuously differentiable in both $F$ and $\theta$ with for some $\delta>0$

$$
E\left[\sup _{\left\|\theta-\theta_{0}\right\| \leq \delta,\left\|F^{\prime}-F\right\| \leq \delta}\left\|\frac{\partial^{j} \psi}{\partial F^{j_{1}} \partial^{j_{2}} \theta}\left(F^{\prime}, Z ; \theta\right)\right\|^{r}\right]<\infty
$$

for $j=0,1,2$ with $j_{1}+j_{2}=j$ and some $r>2$.
B4. $W_{T} \xrightarrow{P} W$, where $W$ is a symmetric positive definite matrix. The matrix $\Gamma_{0}$ is of full rank.
For convenience we assume that $\psi\left(F_{t}, Z_{t} ; \theta_{0}\right)$ is a martingale difference sequence as is plausible in the finance applications we have in mind. Let:

$$
\begin{gathered}
\Gamma_{0}=E\left[\frac{\partial}{\partial \theta} \psi\left(F_{t}, Z_{t} ; \theta_{0}\right)\right] \quad ; \quad V_{0}=\operatorname{var}\left[\psi\left(F_{t}, Z_{t} ; \theta_{0}\right)\right] \\
\Psi=\left(\Gamma_{0}^{\top} W \Gamma_{0}\right)^{-1}\left(\Gamma_{0}^{\top} W V_{0} W \Gamma_{0}\right)\left(\Gamma_{0}^{\top} W \Gamma_{0}\right)^{-1} .
\end{gathered}
$$

Theorem 2. Suppose that Assumptions A1-A7 and B1-B4 hold. Then,

$$
\sqrt{T}\left(\widehat{\theta}-\theta_{0}\right) \Longrightarrow N(0, \Psi)
$$

This shows that the estimation of factors does not affect the limiting distribution of the parameters of the factor process. This means that standard errors for $\widehat{\theta}$ can be constructed as if the factors were observed. Furthermore, under the strong exogeneity assumption of the factors we can factor the likelihood $\ell(y, f \mid X ; g, \theta)=\ell(y \mid X, f ; g) \ell(f \mid X ; \theta)$ so that our two-step approach to estimation of $\theta$ does not lose information. To explore volatility dynamics, we also estimate a vector autoregression using squared factor returns (this model is covered by Theorem 2).

## 5 Empirical Analysis

### 5.1 Data

We follow Fama French (1993) in the construction of the size and value characteristics. For each separate twelve-month period July-June from 1963-2005 we find all securities which have complete CRSP return records over this twelve-month period and the previous twelve month period, and both market capitalization (from CRSP) and book value (from Compustat) records for the previous June. The raw size characteristic each month equals the logarithm of the previous June's market value of equity. The raw value characteristic equals the ratio of the market value of equity to the book value of equity in the previous June. In addition to the Fama-French size and value characteristics we derive from the same return dataset a momentum characteristic as in Carhart (1997). This variable is measured as the cumulative twelve month return up to and including the previous month. Finally we add an own-volatility characteristic, a choice inspired by the recent work of Goyal and Santa Clara (2003) and Ang et al. (2006a,2006b). We define raw volatility as the standard deviation of the individual security return over twelve months up to and including the previous month. The characteristics equal the raw characteristics except standardized each month to have zero mean and unit variance. The size and value characteristics are held constant from July to June whereas the momentum and own-volatility characteristics change each month. Table 1 reports some descriptive statistics for the data: the number of securities in the annual cross section, and the first four crosssectional moments of the four characteristics. To save space the table just shows nine representative dates (July at five-year intervals), as well as time series medians over the full 42 year period, using July data.

Three notes on the interpretation of these characteristics in terms of our econometric theory.

1. We treat all four characteristics as observed without error. Informally, we think of momentum and own-volatility as behaviourally-generated sources of return comovement. Investors observe momentum and own-volatility over the previous twelve months (along with the most recent observations of size and value), adjust their portfolio and pricing behaviour to account for the observed values, and this in turn accounts (for some unspecified reasons) for the subsequent return comovements associated with these characteristics. Understanding more fundamentally the sources of the characteristic-related comovements is an important topic which we do not address here.
2. The cited references Ang et al. (2006a,b) and Goyal and Santa Clara (2003) use idiosyncratic volatility rather than total volatility as a characteristic. From our perspective, total volatility is preferable since it does not require a previous estimation step to remove market-related return from each asset's total return.
3. In our econometric theory we allow all the characteristics to vary freely over time. Since size and value change annually whereas momentum and own-volatility change monthly, another approach would be to modify the econometric theory to allow some characteristics to change only at a lower frequency. We do not pursue this alternative approach here.

A useful descriptive statistic is the correlation matrix of the explanatory variables. This is complicated in our model by the time-varying nature of the characteristics which serve as our explanatory variables. Figure 1 shows for each pair of characteristics the time series evolution of the cross-sectional correlation between them, using the cross-section each July. It is clear that these correlations are not constant over time. The correlation between size and value exhibits slow and persistent swings, with a negative average. Size and momentum on the other hand are on average uncorrelated. Most interesting is the relationship between own-volatility and momentum, taking large swings from high positive correlation of 0.7 to negative correlation of -0.35 . None of the correlations are large enough in magnitude to be worrisome in terms of accurate identification of the model.

### 5.2 Implementation

In the case of a fully balanced panel it would be straightforward to estimate the characteristic-beta function at each data point in the sample. However in the presence of time-varying characteristics this is not feasible since the number of asset returns (each with a unique vector of characteristics) equals $1,886,172$ in our sample. In order to make the algorithm described in Section 3 computationally feasible we concentrate estimation of the characteristic functions on 61 equally-spaced grid points between -3 and 3, which corresponds to a distance of 0.1 between contiguous grid points. We use linear interpolation between the values at these grid points to compute the characteristic-beta function at all $1,886,172$ sample points. Then we use the full sample of $1,886,172$ asset returns and associated factor betas to estimate the factor returns. This procedure greatly improves the speed of our algorithm while sacrificing little accuracy, since the characteristic-beta functions are reasonably linear between these closely-spaced grid points.

We chose a Gaussian kernel throughout to nonparametrically estimate the conditional expectations summarized in (24). The advantage of this kernel is that it is very smooth and produces nice regular estimates, whereas, say the Epanechnikov kernel produces estimates with discontinuities in the second derivatives. The bandwidth choice is done separately for each characteristic function. We follow Connor and Linton (2007) and use their variable bandwidth tied to local data density. For each characteristic value and each year, we calculate the sample density of the root-mean-squared differences between all the sample characteristic and the individual grid point. We then set the bandwidth for this grid point equal to the fifth percentile of this sample density. This implies that
ninety-five percent of the observations are at least one bandwidth away from the grid point, where distance is measured by root-mean-square. This simple procedure guarantees that the bandwidth is narrow where the data set is locally more densely populated (e.g., near the median values of the characteristic) and wider where the data set is locally sparse (e.g., near the extreme values of the characteristic). It is rather like a smooth nearest neighbors bandwidth taking $5 \%$ of the data in each marginal window.

### 5.3 The Characteristic Beta Functions

Table 2 shows the estimates of the characteristic-beta functions at a small selected set of characteristic values and the heteroskedasticity-consistent standard errors from (31) for each of these estimates. To avoid any spurious non-linearity results due to smoothing in regions where there are no data, we report results for each characteristic only over a support ranging from the empirical $2.5 \%$ to the $97.5 \%$ quantile. The standard errors tend to be somewhat larger in the tails, where the data is sparser. Given that our procedure is able to use all 1.8 million return observations to estimate the characteristic-beta functions, the standard errors are small.

The characteristic-beta functions over all grid points are displayed in Figure 2. Note that these characteristic-beta functions satisfy the equally-weighted zero mean/unit variance identification conditions described in Section 3. For comparative purposes we overlay the linear Rosenbergtype model with the same identification conditions imposed. The characteristic-beta functions are mostly monotonically increasing for all four characteristics. Size and value show strongly non-linear characteristic-beta functions, both with concave shapes. The observed shapes for momentum and own-volatility are closer to linear.

### 5.4 Explanatory Power of Each Factor

Note that at each step of the iterative estimation, the factor returns are the coefficients from period-by-period unconstrained cross-sectional regression of returns on the previous iteration's factor betas. To measure the explanatory power of the factors we take the final-step estimates of factor betas and perform the set of cross-sectional regressions with all the factors, each factor singly, and all the factors except each one. Table 3 shows the time-series averages of uncentered $R^{2}$ (UR2) statistic in all these cases: all five factors, each single factor, and each subset of four factors. The market factor is dominant in terms of explanatory power; a well-known result. The own-volatility factor is the strongest of the characteristic-based factors, followed by size, momentum, and value. The ordering of relative importance is the same whether we consider the factors singly or their marginal contribution given the other four.

We test for the statistical significance of each factor by calculating, for each cross-sectional regression, the t-statistic for each estimated coefficient, based on Eicker-White heteroskedasticity-consistent standard errors. Then for each factor we find the average number of cross-sectional regression tstatistics that are significant at a $95 \%$ confidence level across the 504 time periods. The resulting count statistic has an exact binomial distribution under the null hypothesis that the factor return is zero each period. Table 4a shows the annualized means and standard deviations of the factor returns, the percentage of significant t-statistics for each factor, and the aggregate p-value. All five factors are highly significant.

Table 4 b displays the correlations of the estimated factors, along with the three Fama-French factors, RMRF, SMB, and HML. RMRF is the Fama-French market factor, it is the return to the value-weighted market index minus the risk free return; SMB is the return to a small capitalization portfolio minus the return to a large-capitalization portfolio; HML is the return to a high book-toprice portfolio minus the return to a low book-to-price portfolio. See Fama and French (1993) for detailed discussion of their portfolio formation rules. We also include a momentum factor created by Ken French; this is the return to a portfolio with high cumulative returns over the past twelve months minus the return to a portfolio with low cumulative returns over the past twelve months, adjusted to have roughly equal average capitalization; see Ken French's website ${ }^{9}$ for details, where all the Fama-French data is freely available. Our factors and the analogous Fama-French factors are highly correlated (note that the size characteristic is defined inversely in the two models, hence the negative correlation). The Fama-French factors are based on capitalization-weighted portfolios whereas our factors are statistically generated, treating all assets equally. Since the cross-section of securities is dominated, in terms of the number of securities, by low-capitalization firms, this induces a strong positive correlation between our market factor and the Fama-French SMB factor. Our volatility factor has strong positive correlation with the market factor. This corroborates the finding in Ang et al. (2006b) Table 10, which shows high covariance between their idiosyncratic-volatility factor returns and the Fama-French market factor returns. It also seems theoretically consistent with the finding in Ang et al. (2006a) that the market factor return is negatively correlated with changes in VIX, a forward-looking index of market volatility. Essentially, the positive correlation between the own-volatility factor and market factor means that high own-volatility stocks outperform when the overall market rises and underperform when the overall market falls. There is also a strong negative correlation between the own-volatility and momentum factor returns, for which we have no ready explanation.

Figure 3 shows the characteristic-beta functions re-estimated on the four non-overlapping 126month subintervals in the data set. The functions seem stable over time although we do not attempt

[^7]a formal test.

### 5.5 Nonparametric Mispricing Functions

Next we add four additive nonparametric mispricing functions to the model, one for each of the characteristics. The estimation algorithm is essentially unchanged, except that after each iteration we impose the orthogonality and mean-zero identification conditions (6)-(5) on each mispricing function $\widehat{\alpha}_{j}$. We do this by regressing $\widehat{\alpha}_{j}$ on an intercept and $\widehat{g}_{j}$ across the set of grid points, and then replacing $\widehat{\alpha}_{j}$ with the residuals from this regression. The estimated component $\alpha$-functions can be interpreted as a nonparametric version of the tables of estimated alpha coefficients for characteristic-sorted portfolios shown in Fama and French (1993). Figure 4 shows the results. There is little evidence against the five-factor asset pricing model: all four mispricing functions differ only negligibly from zero, mostly lying in the interval from -20 to +20 basis points (in units of monthly return). An obvious feature of all four graphs is the absence of any monotonic pattern upward or downward. On more careful analysis it is clear that this is a consequence of the identification condition (6) rather than an empirical finding. Since the estimated characteristic-beta functions are monotonically increasing, the estimated characteristic-mispricing functions (which are forced to be orthogonal to them) have average slopes close to zero. The same lack-of-monotonicity is apparent in the empirical results of Fama and French (1993); see in particular Table 9a of that paper.

### 5.6 Time Series Dynamics and Trends

Table 5 shows the results from a first-order vector autoregression of the five factors returns on their lagged values. The size factor has the strongest autocorrelation and cross-correlation, as measured by its R-squared in the vector autoregression. This is to be expected, since this factor has the heaviest concentration in low-capitalization, illiquid securities where autocorrelation and cross-correlation is strongest; see, e.g., Lo and Mackinlay (1990). Table 6 shows the first-order vector autoregression using squared factor returns; this formulation is useful for identifying multivariate volatility dynamics in the factor returns. The squared value factor has the highest $R^{2}$ followed by size and own-volatility. It is interesting that the squared market factor shows the lowest $R^{2}$. These results are only intended to be suggestive of the general pattern of dynamic volatility linkages between the factors; to complete the specification would require building a full multivariate volatility model which we do not attempt here, see, e.g., Laurent, Bauwens and Rombouts (2006) and references therein.

Campbell, Lettau, Malkiel and Xu (2001), in a model with industry and market factors, and Jones (2001) and Connor, Korajczyk and Linton (2006) in models with statistical factors, show that
cross-sectional mean-square asset-specific return, $\gamma_{t}=n_{t}^{-1} \sum_{i=1}^{n_{t}} \varepsilon_{i t}^{2}$, varies through time, with a strong upward trend over sample periods covering the 1950s - 1990's. Figure 5 replicates this finding in our model with characteristic-based factors. It is notable that the long upward trend in mean-square asset-specific return seems to reverse in the post-2000 sample period included here. We also examined each of the squared factor returns but found no evidence of similar trends (figures not shown, but available from the authors).

### 5.7 A Weighted Least-Squares Objective Function

Jones (2001) describes how to modify statistical factor estimation methods for the presence of timevarying mean-square asset-specific return as shown in Figure 5. Jones' adjustment is not strictly necessary with our method: although we use a least-squares-type objective function to motivate the estimators, we do not need to assume time-series or cross-sectional homogeneity of asset-specific variances to derive their asymptotic properties. Nonetheless, the evidence in Figure 5 points toward an improvement in the efficiency of the estimators by modifying the objective function (11) to account for the average time-series heteroskedasticity of asset-specific returns. The altered objective function is:

$$
Q_{T}^{*}(f, g)=\frac{1}{T} \sum_{t=1}^{T} \gamma_{t}^{-1} E\left[\left\{y_{i t}-f_{u t}-\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}\right\}^{2}\right]
$$

If $\gamma_{t}$ is the only source of heteroskedasticity in returns, then this adjustment allows us to attain Linton's (1997) oracle bound. Repeating the analysis with this altered objective function the only change to the estimation algorithm is in (15) which becomes:

$$
g_{j}\left(x_{j}\right)=\frac{\frac{1}{T} \sum_{t=1}^{T} \gamma_{t}^{-1} \frac{1}{n_{t}} \sum_{i=1}^{n_{t}} f_{j t}\left(E\left[y_{i t} \mid X_{j i t}=x_{j}\right]-f_{u t}-\sum_{k \neq j} f_{k t} E\left[g_{k}\left(X_{k i t}\right) \mid X_{j i t}=x_{j}\right]\right)}{\frac{1}{T} \sum_{t=1}^{T} \gamma_{t}^{-1} f_{j t}^{2}}
$$

The re-estimated characteristic-beta functions (not shown, available from the authors) differ only very marginally from the original homogenously-time-weighted estimates.

### 5.8 The Case of an Observed Market Factor

Next we re-estimate the model using the market index from Fama and French as $f_{u t}$, and normalizing the characteristic-beta functions using capitalization-weighted means and variances. Figure 6 shows the characteristic-beta functions in this case. The only substantial change from the equally-weighted case is in the location of the characteristic-beta function for size. Now, most firms in the sample have negative size betas rather than, as previously, a roughly equal split between positive and negative size
betas. This reflects the fact that approximately $80 \%$ of the securities have size characteristics below the capitalization-weighted mean characteristic. The own-volatility characteristic-beta function has a steeper positive slope than in the equally-weighted case. Value and Momentum show less of a change from the equally-weighted case.

Table 6 shows tests of the embedded CAPM pricing (8). The well-known size and value premia relative to the CAPM are significant in the overall period, but inconsistent across subperiods. There is a significant negative premium relative to the CAPM for the own-volatility factor. In this fivefactor model, the momentum factor on its own does not have an abnormal return premium relative to the CAPM. The Gibbons-Ross-Shanken joint test is significant in the full period and in three of the four subperiods.

## 6 Summary and Conclusion

Following the pioneering work of Rosenberg (1974), Fama and French (1993) and others, characteristicbased factor models have played a leading role in explaining the comovements of individual equity returns. This paper applies a new weighted additive nonparametric estimation procedure to estimate characteristic-based factor models more data-efficiently than existing nonparametric methods. The methodology we develop extends existing results to the large cross-section long time series setting. We obtain a variety of statistical results that are useful for conducting inference. We think this methodology can be useful elsewhere.

We estimate a characteristic-based factor model with five factors: a market factor, size factor, value factor, momentum factor and own-volatility factor. Although much of the existing literature has focused on the three-factor Fama-French model (market, size and value) we find that the momentum and own-volatility factors are at least, if not more, important than size and value in explaining return comovements. The univariate functions mapping characteristics to factor betas are monotonic but not linear, the deviation from linearity is particularly strong for size and value, less so for momentum and own-volatility.

Our methodology provides a new nonparametric test of multi-beta asset pricing theory. We estimate nonparametrically a set of additive mispricing functions based on the four security characteristics. We find little evidence against the five-factor asset pricing model.

We also examine the time-series behaviour of the estimated factor returns and their squared values using vector autoregressions. There are strong, multivariate dynamics in the volatilities of the factor returns.

We consider the case in which the market factor is observed (set equal to a capitalization-weighted index) and test an embedded CAPM pricing relationship against the more general five-factor pricing
model. We reject the CAPM restriction relative to the five-factor alternative, due to abnormal return premia associated with the size, value and own-volatility factors.

The model we examine can be extended in various ways. One may want to combine features of the statistical factor model (2) with our observed factor model (1) by specifying

$$
\begin{equation*}
y_{i t}=f_{u t}+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}+\sum_{k=J+1}^{K} \beta_{k i} f_{k t}+\varepsilon_{i t} . \tag{34}
\end{equation*}
$$

Our analysis can be extended to this case: first, apply our procedures above as if the unobserved components were not there, and then apply methods appropriate for statistical factor analysis to the residuals. Under some conditions on the time series dependence of $f_{t}$, our results should carry over to this case. ${ }^{10}$ In other applications one might think there are observable time series variables $Z_{j t}$ that affect the common components $f_{j t}$, e.g., $f_{j t}=h_{j}\left(Z_{j t}\right)$ for some functions $h_{j}$. This case can also be analyzed in a similar way.

## A Appendix

## A. 1 Population Integral Equation

Here we establish the properties of the population equations (14). Define

$$
\begin{gathered}
m_{j}^{T}(x)=\frac{\sum_{t=1}^{T} f_{j t} E\left[\left(y_{i t}-f_{u t}\right) \mid X_{j i}=x\right]}{\sum_{t=1}^{T} f_{j t}^{2}} \\
\mathcal{H}_{j k}^{T}\left(x, x^{\prime}\right)=\rho_{j k}^{T} \times \frac{p_{k, j}\left(x, x^{\prime}\right)}{p_{j}(x) p_{k}\left(x^{\prime}\right)}, \text { where } \rho_{j k}^{T}=\frac{\sum_{t=1}^{T} f_{j t} f_{k t}}{\sum_{t=1}^{T} f_{j t}^{2}}
\end{gathered}
$$

and $p_{k, j}$ is the joint density of $\left(X_{j i}, X_{k i}\right)$. We drop the $T$ superscript in $\mathcal{H}_{j k}^{T}, m_{j}^{T}$, and $\rho_{j k}^{T}$ for the remainder of this subsection.

The equations (14) can be rewritten as

$$
\left(\begin{array}{cccc}
I & \mathcal{H}_{1} & \cdots & \mathcal{H}_{1}  \tag{35}\\
\mathcal{H}_{2} & I & \mathcal{H}_{2} & \cdots \\
\vdots & & \ddots & \\
\mathcal{H}_{J} & \cdots & & I
\end{array}\right)\left(\begin{array}{c}
g_{1} \\
g_{2} \\
\vdots \\
g_{J}
\end{array}\right)=\mathcal{H} g=m=\left(\begin{array}{c}
m_{1} \\
m_{2} \\
\vdots \\
m_{J}
\end{array}\right)
$$

[^8]where $I$ is the identity operator and $\left(\mathcal{H}_{j} g_{k}\right)(x)=\int \mathcal{H}_{j k}\left(x, x^{\prime}\right) g_{k}\left(x^{\prime}\right) p_{k}\left(x^{\prime}\right) d x^{\prime}$. Note that (35) is a system of linear type 2 integral equations in the functions $g_{j}$. It is very similar to the first order conditions associated with additive nonparametric regression with two exceptions. First, in that case the intercept function is just an unweighted conditional expectation $E\left[y_{i} \mid X_{j i}=x\right]$. Second, the operator in that case does not have the weighting $\rho_{j k}$. The first difference is irrelevant for the studying of the existence and uniqueness of solutions, since only the operator is required for that. The second difference is rather minor since the weighting factors $\rho_{j k}$ do not vary with the covariates and by Assumption A5 the $\rho_{j k}$ are uniformly bounded, since
\[

$$
\begin{equation*}
\left(\sum_{t=1}^{T} f_{j t} f_{k t}\right)^{2} \leq\left(\sum_{t=1}^{T} f_{j t}^{2}\right)\left(\sum_{t=1}^{T} f_{k t}^{2}\right) \tag{36}
\end{equation*}
$$

\]

by the Cauchy Schwarz inequality. Suppose also that the Hilbert-Schmidt condition holds:

$$
\begin{equation*}
\int \frac{p_{k, j}\left(x, x^{\prime}\right)^{2}}{p_{j}(x) p_{k}\left(x^{\prime}\right)} d x d x^{\prime}<\infty, \quad \text { for all } j, k \tag{37}
\end{equation*}
$$

This is satisfied under our assumption A2. Then let $\Psi_{j}$ be the operator such that $\Psi_{j} f_{j}\left(x_{j}\right)=0$ and $\Psi_{j} f_{k}\left(x_{k}\right)=f_{k}\left(x_{k}\right)-\int \mathcal{H}_{j k}\left(x_{k}, u\right) f_{k}(u) p_{k}(u) d u$ for any functions $f_{j}, f_{k}$, and define $\mathcal{T}=\Psi_{J} \cdots \Psi_{1}$. This operator represents a population version of one cycle of the iteration (17). It can be shown from Lemma 1 of Mammen, Linton, and Nielsen (1999) that $\mathcal{T}$ is a positive self-adjoint linear operator with operator norm less than one, which implies that the population iterations converge to the solution of (35) at a geometric rate. Consider the special case of $J=2$, in which case it suffices to show that the operator norm of $\mathcal{H}_{1} \mathcal{H}_{2}$ is strictly less than one. Consider the classical bivariate additive regression on covariates $X_{1}, X_{2}$, with associated operators $\widetilde{\mathcal{H}}_{1}$ and $\widetilde{\mathcal{H}}_{2}$, Mammen, Linton, and Nielsen (1999). Then $\mathcal{H}_{1} \mathcal{H}_{2}=\rho_{12} \rho_{21} \widetilde{\mathcal{H}}_{1} \widetilde{\mathcal{H}}_{2}$, where

$$
\rho_{12} \rho_{21}=\frac{\left(\sum_{t=1}^{T} f_{1 t} f_{2 t}\right)^{2}}{\sum_{t=1}^{T} f_{1 t}^{2} \sum_{t=1}^{T} f_{2 t}^{2}} \leq 1
$$

by (36). It follows that the operator norm of $\mathcal{H}_{1} \mathcal{H}_{2}$ is strictly less than one if and only if the operator norm of $\widetilde{\mathcal{H}}_{1} \widetilde{\mathcal{H}}_{2}$ is strictly less than one. This is implied by the Hilbert-Schmidt condition (37).

## A. 2 Proof of Results

Proof of Theorem 1. The proof strategy is to first establish the properties of the initial consistent estimators and then to work with iterations from these starting points. As it turns out, this strategy obviates the need to establish the convergence of the algorithm and to deal with the integral equation (35) in any detail. Take the time averaged data (20) and estimate the functions $\bar{g}_{j}($.$) by the smooth$
backfitting method and then renormalize. The properties of the resulting estimator $\widetilde{g}_{j}\left(x_{j}\right), j=$ $1, \ldots, J$, are as in Mammen, Linton, and Nielsen (1999) except that the error term is $O_{p}\left(T^{-1 / 2}\right)$. In particular, we have for an interior point $x_{j}$

$$
\begin{equation*}
\widetilde{g}_{j}\left(x_{j}\right)-g_{j}\left(x_{j}\right)=\frac{1}{n p_{j}\left(x_{j}\right)} \sum_{i=1}^{n} K_{h}\left(x_{j}-X_{j i}\right) \bar{\varepsilon}_{i}+h^{2} \beta_{n, j}\left(x_{j}\right)+\frac{1}{n} \sum_{i=1}^{n} s_{n}\left(X_{i}, x_{j}\right) \bar{\varepsilon}_{i}+R_{n j}\left(x_{j}\right), \tag{38}
\end{equation*}
$$

where $\sup _{x_{j}}\left|R_{n j}\left(x_{j}\right)\right|=o_{p}\left(n^{-1 / 2} T^{-1 / 2}\right), \beta_{n, j}\left(x_{j}\right)$ is stochastically bounded with $\sup _{x_{j}} \mid \beta_{n, j}\left(x_{j}\right)-$ $\beta_{j}\left(x_{j}\right) \mid=o_{p}(1)$, where $\beta_{j}\left(x_{j}\right)$ is a deterministic bounded continuous function. The function $s_{n}$ is stochastically bounded (and depends only on $\left.X_{1}, \ldots, X_{n}\right)$ with $E s_{n}\left(X_{i}, x_{j}\right)^{2}<\infty$. The error term $\bar{\varepsilon}_{i}$ has variance of order $T^{-1}$, which implies that the leading term in the expansion is $O_{p}\left(n^{-1 / 2} h^{-1 / 2} T^{-1 / 2}\right)+$ $O_{p}\left(h^{2}\right)$. Note that the remainder contains terms of the form $O_{p}\left(n^{-1 / 2} h^{-1 / 2}\right) \times O_{p}\left(n^{-1 / 2} h^{-1 / 2} T^{-1 / 2}\right)$, which are $o_{p}\left(n^{-1 / 2} T^{-1 / 2}\right)$ under our bandwidth conditions. It follows that $\widetilde{g}_{j}\left(x_{j}\right)$ is $\sqrt{n T h}$ consistent and asymptotically normal, and asymptotically independent of $\widetilde{g}_{k}\left(x_{k}\right)$. Note that the renormalization only affects the bias and not the variance due to the effect that integration has on variance.

The proof of our main result is given in the following lemmas. Lemma 1 and 2 give the pointwise performance of the initial factor estimator, denoted $\widetilde{f}_{t}$, while Lemmas 3 and 4 give the uniform (over $t)$ performance of $\widetilde{f}_{t}$. Lemma 5 and 6 give the pointwise expansion of the one-step update $\widetilde{g}_{j}^{[1]}\left(x_{j}\right)$ of $\widetilde{g}_{j}\left(x_{j}\right)$. We then generalize to the $k$-step update covered by the theorem.

Consider the infeasible estimator $f_{t}^{\dagger}$ that is any solution of the system of linear equations $A^{\dagger} f_{t}^{\dagger}=$ $b_{t}^{\dagger}$, where

$$
b_{t}^{\dagger}=\frac{1}{n} \sum_{i=1}^{n} G\left(X_{i}\right) y_{i t} \quad ; \quad A^{\dagger}=\frac{1}{n} \sum_{i=1}^{n} G\left(X_{i}\right) G\left(X_{i}\right)^{\top} .
$$

This is just a standard OLS estimator with regressors one and $g_{j}\left(X_{j i}\right)$. Let also $A=E\left[G\left(X_{i}\right) G\left(X_{i}\right)^{\top}\right]$.
Lemma 1. Under our assumptions, for any $t$

$$
\sqrt{n}\left(f_{t}^{\dagger}-f_{t}\right) \Longrightarrow N\left(0, V_{t, t}\right)
$$

and $f_{t}^{\dagger}, f_{s}^{\dagger}$ are asymptotically independent for any $s \neq t$.
Proof. We have

$$
\sqrt{n}\left[\begin{array}{l}
f_{t}^{\dagger}-f_{t} \\
f_{s}^{\dagger}-f_{s}
\end{array}\right]=\left(A^{\dagger}\right)^{-1}\binom{\frac{1}{\sqrt{n}} \sum_{i=1}^{n} G\left(X_{i}\right) \varepsilon_{i t}}{\frac{1}{\sqrt{n}} \sum_{i=1}^{n} G\left(X_{i}\right) \varepsilon_{i s}} .
$$

Then by the cross-sectional independence, $A^{\dagger}=A+o_{p}(1)$. Let

$$
u_{i, t s}(c)=c^{\top}\left[\begin{array}{l}
G\left(X_{i}\right) \varepsilon_{i t} \\
G\left(X_{i}\right) \varepsilon_{i s}
\end{array}\right]
$$

for some vector $c$. Then by the Lindeberg CLT,

$$
\frac{1}{\sqrt{n}} \sum_{i=1}^{n} u_{i, t s}(c) \Longrightarrow N\left(0, \sigma_{u, t s}^{2}(c)\right),
$$

where $\sigma_{u, t s}^{2}(c)=\operatorname{var}\left[u_{i, t s}(c)\right]=c^{\top} E\left[\varepsilon_{i t}^{2} G\left(X_{i}\right) G\left(X_{i}\right)^{\top}\right] c<\infty$. The result then follows by the CramerRao device, Slutsky theorem, and the fact that $\operatorname{cov}\left(\varepsilon_{i t}, \varepsilon_{i s}\right)=0$.

Now consider the feasible factor estimator based on the time-averaged backfitting estimator $\widetilde{f_{t}}=$ $\widetilde{A}^{-1} \widetilde{b}_{t}$, where

$$
\widetilde{b}_{t}=\frac{1}{n} \sum_{i=1}^{n} \widetilde{G}\left(X_{i}\right) y_{i t} \quad ; \quad \widetilde{A}=\frac{1}{n} \sum_{i=1}^{n} \widetilde{G}\left(X_{i}\right) \widetilde{G}\left(X_{i}\right)^{\top},
$$

where $\widetilde{G}\left(X_{i}\right)=\left[1, \widetilde{g}_{1}\left(X_{1 i}\right), \ldots, \widetilde{g}_{J}\left(X_{J i}\right)\right]^{\top}$.
Lemma 2. Under our assumptions, for any $t$ there is a stochastically bounded sequence $\delta_{n, t}$ such that

$$
\sqrt{n}\left(\widetilde{f_{t}}-f_{t}^{\dagger}-h^{2} \delta_{n, t}\right)=o_{p}(1) .
$$

Proof. We use the matrix expansion $(I+\Delta)^{-1}=I-\Delta+(I+\Delta)^{-1} \Delta^{2}$ to obtain

$$
\begin{align*}
\widetilde{f}_{t}-f_{t}^{\dagger}= & \widetilde{A}^{-1} \widetilde{b}_{t}-A^{\dagger-1} b_{t}^{\dagger} \\
= & A^{\dagger-1}\left(\widetilde{b}_{t}-b_{t}^{\dagger}\right)-A^{\dagger-1}\left(\widetilde{A}-A^{\dagger}\right) A^{\dagger-1} b_{t}^{\dagger}  \tag{39}\\
& -A^{\dagger-1}\left(\widetilde{A}-A^{\dagger}\right) A^{\dagger-1}\left(\widetilde{b}_{t}-b_{t}^{\dagger}\right) \\
& +A^{\dagger-1 / 2}\left[I+A^{\dagger-1 / 2}\left(\widetilde{A}-A^{\dagger}\right) A^{\dagger-1 / 2}\right]^{-1} \widetilde{A}^{-1 / 2}\left(\widetilde{A}-A^{\dagger}\right) A^{\dagger-1}\left(\widetilde{A}-A^{\dagger}\right) A^{\dagger-1} \widetilde{b}_{t} .
\end{align*}
$$

The error $\left\|\widetilde{f}_{t}-f_{t}^{\dagger}\right\|$ is majorized by the errors $\left\|\widetilde{b}_{t}-b_{t}^{\dagger}\right\|$ and $\left\|\widetilde{A}-A^{\dagger}\right\|$ times constants due to the invertibility of $A^{\dagger}$. For example,

$$
\begin{gathered}
\left\|A^{\dagger-1}\left(\widetilde{b}_{t}-b_{t}^{\dagger}\right)\right\| \leq \lambda_{\max }\left(A^{\dagger-1}\right)\left\|\widetilde{b}_{t}-b_{t}^{\dagger}\right\|=\frac{1}{\lambda_{\min }\left(A^{\dagger}\right)}\left(\sum_{j=0}^{J}\left(\widetilde{b}_{j t}-b_{j t}^{\dagger}\right)^{2}\right)^{1 / 2}, \\
\left\|A^{\dagger-1}\left(\widetilde{A}-A^{\dagger}\right) A^{\dagger-1} b_{t}^{\dagger}\right\| \leq \frac{\lambda_{\max }^{1 / 2}\left(\left(\left(\widetilde{A}-A^{\dagger}\right)^{\top}\left(\widetilde{A}-A^{\dagger}\right)\right)\right.}{\lambda_{\min }^{2}\left(A^{\dagger}\right)}\left\|b_{t}^{\dagger}\right\| \leq \frac{\left(\sum_{j, k=0}^{J}\left(\widetilde{A}_{j k}-A_{j k}^{\dagger}\right)^{2}\right)^{1 / 2}}{\lambda_{\min }^{2}\left(A^{\dagger}\right)}\left\|b_{t}^{\dagger}\right\|,
\end{gathered}
$$

where $\lambda_{\max }($.$) and \lambda_{\min }($.$) denote the largest and smallest (respectively) eigenvalues of a square$ symmetric matrix. Furthermore, $\lambda_{\min }\left(A^{\dagger}\right) \geq \lambda_{\min }(A)-o_{p}(1)$, where by assumption, $\lambda_{\min }(A)>0$. We establish the order in probability of the terms $\widetilde{b}_{j t}-b_{j t}^{\dagger}$ and $\widetilde{A}_{j k}-A_{j k}^{\dagger}$.

Consider the typical element in $\widetilde{b}_{t}-b_{t}^{\dagger}$,

$$
\begin{aligned}
\frac{1}{n} \sum_{i=1}^{n} y_{i t}\left[\widetilde{g}_{j}\left(X_{j i}\right)-g_{j}\left(X_{j i}\right)\right]= & \frac{1}{n} \sum_{i=1}^{n}\left[f_{u t}+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}\right]\left[\widetilde{g}_{j}\left(X_{j i}\right)-g_{j}\left(X_{j i}\right)\right] \\
& +\frac{1}{n} \sum_{i=1}^{n} \varepsilon_{i t}\left[\widetilde{g}_{j}\left(X_{j i}\right)-g_{j}\left(X_{j i}\right)\right] \\
= & T_{n 1}+T_{n 2} .
\end{aligned}
$$

We consider the term $T_{n 1}$. From (38), we have

$$
\begin{aligned}
T_{n 1}= & \frac{1}{n} \sum_{i=1}^{n}\left[f_{u t}+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}\right]\left[\widetilde{g}_{j}\left(X_{j i}\right)-g_{j}\left(X_{j i}\right)\right] \\
= & \frac{1}{n} \sum_{i=1}^{n}\left[f_{u t}+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}\right] \frac{1}{n p_{j}\left(X_{j i}\right)} \sum_{i^{\prime}=1}^{n} K_{h}\left(X_{j i}-X_{j i^{\prime}}\right) \bar{\varepsilon}_{i^{\prime}} \\
& +h^{2} \frac{1}{n} \sum_{i=1}^{n}\left[f_{u t}+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}\right] \beta_{n, j}\left(X_{j i}\right) \\
& +\frac{1}{n} \sum_{i=1}^{n}\left[f_{u t}+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}\right] \frac{1}{n} \sum_{i^{\prime}=1}^{n} s_{n}\left(X_{i^{\prime}}, x_{j}\right) \bar{\varepsilon}_{i^{\prime}} \\
& +\frac{1}{n} \sum_{i=1}^{n}\left[f_{u t}+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}\right] R_{n j}\left(X_{j i}\right) \\
= & T_{n 11}+T_{n 12}+T_{n 13}+T_{n 14} .
\end{aligned}
$$

The first term $T_{n 11}$ is a degenerate U-statistic, Powell, Stock, and Stoker (1989), i.e., $T_{n 11}=$ $\sum \sum_{i, i^{\prime}} \varphi_{n i, i^{\prime}}$ with

$$
\varphi_{n i, i^{\prime}}=\frac{1}{n^{2}}\left[f_{u t}+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}\right] \frac{K_{h}\left(X_{j i}-X_{j i^{\prime}}\right)}{p_{j}\left(X_{j i}\right)} \bar{\varepsilon}_{i^{\prime}}
$$

where $E\left[\varphi_{n i, i^{\prime}} \mid X_{i}, \bar{\varepsilon}_{i}\right]=0$. Therefore, we can write

$$
T_{n 11}=\sum_{i=1}^{n} \varphi_{n i, i}+\sum_{i^{\prime}=1}^{n} E\left[\varphi_{n i, i^{\prime}} \mid X_{i^{\prime}}, \bar{\varepsilon}_{i^{\prime}}\right]+\sum_{i \neq i^{\prime}} \sum_{n i, i^{\prime}},
$$

where $\widetilde{\varphi}_{n i, i^{\prime}}=\varphi_{n i, i^{\prime}}-E\left[\varphi_{n i, i^{\prime}} \mid X_{i^{\prime}}, \overline{\bar{i}}_{i^{\prime}}\right]$ and by construction $E\left[\widetilde{\varphi}_{n i, i^{\prime}} \mid X_{i}, \bar{\varepsilon}_{i}\right]=E\left[\widetilde{\varphi}_{n i, i^{\prime}} \mid X_{i^{\prime}}, \bar{\varepsilon}_{i^{\prime}}\right]=0$. By straightforward moment calculations it can be shown that $\sum_{i=1}^{n} \varphi_{n i, i}=O_{p}\left(n^{-3 / 2} h^{-1} T^{-1 / 2}\right)=$ $o_{p}\left((n T)^{-1 / 2}\right)$. Furthermore, $\operatorname{var}\left(\widetilde{\varphi}_{n i, i^{\prime}}\right)=O\left(n^{-4} T^{-1} h^{-1}\right)$ so that $\sum \sum_{i \neq i^{\prime}} \widetilde{\varphi}_{n i, i^{\prime}}=O_{p}\left(n^{-1} T^{-1 / 2} h^{-1 / 2}\right)=$ $o_{p}\left((n T)^{-1 / 2}\right)$. Finally,

$$
\sum_{i^{\prime}=1}^{n} E\left[\varphi_{n i, i^{\prime}} \mid X_{i^{\prime}}, \bar{\varepsilon}_{i^{\prime}}\right] \simeq \frac{1}{n} \sum_{i=1}^{n} \bar{\varepsilon}_{i}\left[f_{u t}+\sum_{j^{\prime}=1}^{J} E\left[g_{j^{\prime}}\left(X_{j^{\prime} i}\right) \mid X_{j i}\right] f_{j^{\prime} t}\right]=O_{p}\left((n T)^{-1 / 2}\right)
$$

Furthermore, $T_{n 12}=O_{p}\left(h^{2}\right)$ and $T_{n 13}, T_{n 14}=o_{p}\left((n T)^{-1 / 2}\right)$. Therefore, $\widetilde{b}_{t}-b_{t}^{\dagger}=O_{p}\left((n T)^{-1 / 2}\right)+O_{p}\left(h^{2}\right)$.
Likewise the typical element in $\widetilde{A}-A^{\dagger}$ satisfies

$$
\frac{1}{n} \sum_{i=1}^{n}\left[\widetilde{g}_{j}\left(X_{j i}\right) \widetilde{g}_{k}\left(X_{k i}\right)-g_{j}\left(X_{j i}\right) g_{k}\left(X_{k i}\right)\right]=O_{p}\left(h^{2}\right)+O_{p}\left((n T)^{-1 / 2}\right) .
$$

It follows that provided $n h^{4} \rightarrow 0, \sqrt{n}\left(\widetilde{f_{t}}-f_{t}^{\dagger}\right)=o_{p}(1)$. More generally, letting

$$
\delta_{n, t}=\frac{1}{n} \sum_{i=1}^{n}\left[f_{u t}+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}\right] \beta_{n, j}\left(X_{j i}\right)
$$

we have $\sqrt{n}\left(\widetilde{f}_{t}-f_{t}^{\dagger}-h^{2} \delta_{n, t}\right)=o_{p}(1)$.
We now turn to the uniform over $t$ properties, (28). By the triangle inequality

$$
\max _{1 \leq t \leq T}\left\|\widetilde{f}_{t}-f_{t}\right\| \leq \max _{1 \leq t \leq T}\left\|\widetilde{f}_{t}-f_{t}^{\dagger}\right\|+\max _{1 \leq t \leq T}\left\|f_{t}^{\dagger}-f_{t}\right\| .
$$

We first examine $\max _{1 \leq t \leq T}\left\|f_{t}^{\dagger}-f_{t}\right\|$.
Lemma 3. Under our assumptions

$$
\max _{1 \leq t \leq T}\left\|f_{t}^{\dagger}-f_{t}\right\|=O_{p}\left(n^{-1 / 2}(\log T)^{\rho}\right)
$$

Proof. By the above arguments, there is a finite constant $C$ such that

$$
\max _{1 \leq t \leq T}\left\|f_{t}^{\dagger}-f_{t}\right\| \leq\left(C+o_{p}(1)\right) \max _{j=u, 1, \ldots, J} \max _{1 \leq t \leq T}\left|\frac{1}{n} \sum_{i=1}^{n} \varepsilon_{i t} g_{j}\left(X_{j i}\right)\right| .
$$

So it suffices to show that

$$
\begin{equation*}
\max _{1 \leq t \leq T}\left|\frac{1}{n} \sum_{i=1}^{n} \varepsilon_{i t} g_{j}\left(X_{j i}\right)\right|=O_{p}\left(n^{-1 / 2}(\log T)^{\rho}\right) \tag{40}
\end{equation*}
$$

for some $\rho>0$. Let $\varepsilon_{i t}^{+}=\varepsilon_{i t} 1\left(\left|\varepsilon_{i t}\right| \leq(n T)^{1 / \kappa}\right)-E\left[\varepsilon_{i t} 1\left(\left|\varepsilon_{i t}\right| \leq(n T)^{1 / \kappa}\right)\right]$. Then, $1-\operatorname{Pr}\left[\left|\varepsilon_{i t}\right| \leq(n T)^{1 / \kappa}\right.$ for $1 \leq t \leq T$ and $1 \leq i \leq n] \leq n T \operatorname{Pr}\left[\left|\varepsilon_{i t}\right|>(n T)^{1 / \kappa}\right] \leq E\left[\left|\varepsilon_{i t}\right|^{\kappa} 1\left(\left|\varepsilon_{i t}\right|>(n T)^{1 / \kappa}\right)\right] \rightarrow 0$. We now apply the Bonferroni and exponential inequalities to $\max _{1 \leq t \leq T}\left|\frac{1}{n} \sum_{i=1}^{n} \varepsilon_{i t}^{+} g_{j}\left(X_{j i}\right)\right|$. In particular, letting $\tau_{n T}^{2}=\inf _{1 \leq t \leq T} \operatorname{var}\left[\sum_{i=1}^{n} \varepsilon_{i t}^{+} g_{j}\left(X_{j i}\right)\right]$, we have

$$
\begin{aligned}
\operatorname{Pr}\left[\max _{1 \leq t \leq T}\left|\sum_{i=1}^{n} \varepsilon_{i t}^{+} g_{j}\left(X_{j i}\right)\right|>K n^{1 / 2}\right] & \leq \sum_{t=1}^{T} \operatorname{Pr}\left[\left|\sum_{i=1}^{n} \varepsilon_{i t}^{+} g_{j}\left(X_{j i}\right)\right|>K n^{1 / 2}\right] \\
& \leq 2 T \exp \left(-\frac{n K^{2}}{2 \tau_{n T}^{2}+2(n T)^{1 / \kappa} n^{1 / 2} K / 3}\right) .
\end{aligned}
$$

By taking $K=(\log T)^{\rho}$ the right hand side is $o(1)$ provided $\kappa>4$.

## Lemma 4. Under our assumptions

$$
\max _{1 \leq t \leq T}\left\|\widetilde{f}_{t}-f_{t}^{\dagger}\right\|=O_{p}\left((n T)^{-1 / 2}(\log T)^{\rho}\right)+O_{p}\left(h^{2}(\log T)^{\rho \prime}\right) .
$$

Proof. As before we apply the triangle inequality again to each term in (39). We have $\max _{1 \leq t \leq T}\left\|b_{t}^{\dagger}\right\|=O_{p}\left((\log T)^{\rho \prime}\right)$, so it suffices to bound the terms $\max _{1 \leq t \leq T}\left|\widetilde{b}_{j t}-b_{j t}^{\dagger}\right|$ and $\max _{1 \leq t \leq T} \mid \widetilde{A}_{j k}-$ $A_{j k}^{\dagger} \mid$. We just show that

$$
\begin{align*}
\max _{1 \leq t \leq T}\left|\frac{1}{n} \sum_{i=1}^{n}\left[f_{u t}+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}\right]\left[\widetilde{g}_{j}\left(X_{j i}\right)-g_{j}\left(X_{j i}\right)\right]\right| & =O_{p}\left(n^{-1 / 2}(\log T)^{\rho}\right)+O_{p}\left(h^{2}(\log T)^{\rho \prime}\right)  \tag{41}\\
\max _{1 \leq t \leq T}\left|\frac{1}{n} \sum_{i=1}^{n} \varepsilon_{i t}\left(\widetilde{g}_{j}\left(X_{j i}\right)-g_{j}\left(X_{j i}\right)\right]\right| & =O_{p}\left(n^{-1 / 2}(\log T)^{\rho}\right) . \tag{42}
\end{align*}
$$

This uses the same type of techniques as above. In particular, we have

$$
\begin{aligned}
& \max _{1 \leq t \leq T}\left|h^{2} \frac{1}{n} \sum_{i=1}^{n}\left[f_{u t}+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}\right] \beta_{j}\left(X_{j i}\right)\right| \\
\leq & h^{2}\left(\max _{1 \leq t \leq T}\left|f_{u t}\right| \frac{1}{n} \sum_{i=1}^{n}\left|\beta_{j}\left(X_{j i}\right)\right|+\sum_{j=1}^{J} \max _{1 \leq t \leq T}\left|f_{j t}\right| \frac{1}{n} \sum_{i=1}^{n}\left|g_{j}\left(X_{j i}\right)\right|\left|\beta_{j}\left(X_{j i}\right)\right|\right) \\
= & O_{p}\left(h^{2}(\log T)^{\rho^{\prime}}\right) .
\end{aligned}
$$

Furthermore,

$$
\max _{1 \leq t \leq T}\left|\frac{1}{n} \sum_{i=1}^{n} \bar{\varepsilon}_{i}\left[f_{u t}+\sum_{j^{\prime}=1}^{J} E\left[g_{j^{\prime}}\left(X_{j^{\prime} i}\right) \mid X_{j i}\right] f_{j^{\prime} t}\right]\right|=O_{p}\left((n T)^{-1 / 2}(\log T)^{\rho^{\prime}}\right) .
$$

In conclusion we have shown $\max _{1 \leq t \leq T}\left\|\widetilde{f}_{t}-f_{t}\right\|=O_{p}\left((n T)^{-1 / 2}(\log T)^{\rho}\right)+O_{p}\left(h^{2}(\log T)^{\rho^{\prime}}\right)$.
Finally, we establish the asymptotic distribution of $\widehat{g}_{j}\left(x_{j}\right)$. Consider the one-step estimator

$$
\widehat{g}_{j}^{[1]}\left(x_{j}\right)=\frac{\sum_{t=1}^{T} \tilde{f}_{j t}\left[\widehat{\lambda}_{1 t}\left(j, x_{j}\right)-\tilde{f}_{u t}-\sum_{k \neq j} \tilde{f}_{k t} \tilde{\lambda}_{2}\left(j, k, x_{j}\right)\right]}{\sum_{t=1}^{T} \tilde{f}_{j t}^{2}},
$$

where

$$
\widetilde{\lambda}_{2}\left(j, k, x_{j}\right)=\frac{\sum_{i=1}^{n} K_{h}\left(X_{j i}-x_{j}\right) \widetilde{g}_{k}\left(X_{k i}\right)}{\sum_{i=1}^{n} K_{h}\left(X_{j i}-x_{j}\right)} .
$$

Define

$$
\widetilde{g}_{j}^{[1]}\left(x_{j}\right)=\frac{\sum_{t=1}^{T} f_{j t}\left[\widehat{\lambda}_{1 t}\left(j, x_{j}\right)-f_{u t}-\sum_{k \neq j} f_{k t} \widetilde{\lambda}_{2}\left(j, k, x_{j}\right)\right]}{\sum_{t=1}^{T} f_{j t}^{2}} .
$$

Lemma 5. Under our assumptions

$$
\widehat{g}_{j}^{[1]}\left(x_{j}\right)-g_{j}\left(x_{j}\right)=\widetilde{g}_{j}^{[1]}\left(x_{j}\right)-g_{j}\left(x_{j}\right)+O_{p}\left((n T)^{-1 / 2}\right)+O_{p}\left(h^{2}\right) .
$$

Proof. We expand out $\widehat{g}_{j}^{[1]}\left(x_{j}\right)$ about $\widetilde{g}_{j}^{[1]}\left(x_{j}\right)$ in a Taylor expansion in $\widetilde{f}_{j t}-f_{j t}$ and $\widetilde{g}_{k}\left(X_{k i}\right)-$ $g_{k}\left(X_{k i}\right)$ obtaining many terms. A typical term is

$$
\frac{\sum_{t=1}^{T}\left(\widetilde{f}_{j t}-f_{j t}\right) f_{j t} g_{j}\left(x_{j}\right)}{\sum_{t=1}^{T} f_{j t}^{2}}
$$

Then

$$
\begin{equation*}
\frac{1}{T} \sum_{t=1}^{T}\left(\widetilde{f}_{j t}-f_{j t}\right) f_{j t}=\frac{1}{T} \sum_{t=1}^{T}\left(\widetilde{f}_{j t}-f_{j t}^{\dagger}\right) f_{j t}+\frac{1}{T} \sum_{t=1}^{T}\left(f_{j t}^{\dagger}-f_{j t}\right) f_{j t} \tag{43}
\end{equation*}
$$

where

$$
\frac{1}{T} \sum_{t=1}^{T}\left(f_{j t}^{\dagger}-f_{j t}\right) f_{j t}=\left(A^{\dagger}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} \frac{1}{n} \sum_{i=1}^{n} G\left(X_{i}\right) \varepsilon_{i t}\left(f_{j t}^{\dagger}-f_{j t}\right) f_{j t}=O_{p}\left((n T)^{-1 / 2}\right)
$$

The expansion for $T^{-1} \sum_{t=1}^{T}\left(\widetilde{f}_{j t}-f_{j t}^{\dagger}\right) f_{j t}$ is more complicated but basically one obtains terms like

$$
\begin{aligned}
& \frac{1}{T} \sum_{t=1}^{T} f_{j t} \frac{1}{n} \sum_{i=1}^{n} y_{i t}\left[\widetilde{g}_{j}\left(X_{j i}\right)-g_{j}\left(X_{j i}\right)\right] \\
= & \frac{1}{T} \sum_{t=1}^{T} f_{j t} \frac{1}{n} \sum_{i=1}^{n} \varepsilon_{i t}\left[\widetilde{g}_{j}\left(X_{j i}\right)-g_{j}\left(X_{j i}\right)\right] \\
& +\frac{1}{T} \sum_{t=1}^{T} f_{j t} \frac{1}{n} \sum_{i=1}^{n}\left[f_{u t}+\sum_{j=1}^{J} g_{j}\left(X_{j i}\right) f_{j t}\right]\left[\widetilde{g}_{j}\left(X_{j i}\right)-g_{j}\left(X_{j i}\right)\right] .
\end{aligned}
$$

Substituting in the expansion (38) and collecting terms we see that the double averaging makes the stochastic terms $O_{p}\left((n T)^{-1 / 2}\right)$. The bias terms are always $O_{p}\left(h^{2}\right)$.

Define

$$
\begin{gather*}
\widetilde{U}_{n 1}=\frac{1}{p_{j}\left(x_{j}\right)} \frac{1}{n} \sum_{i=1}^{n} K_{h}\left(x_{j}-X_{j i}\right) \widetilde{\varepsilon}_{i},  \tag{44}\\
\widetilde{\varepsilon}_{i}=\frac{\sum_{t=1}^{T} f_{j t} \varepsilon_{i t}}{\sum_{t=1}^{T} f_{j t}^{2}} .
\end{gather*}
$$

Lemma 6. Under our assumptions

$$
\widetilde{g}_{j}^{[1]}\left(x_{j}\right)-g_{j}\left(x_{j}\right)=\widetilde{U}_{n 1}+O_{p}\left(h^{2}\right)+o_{p}\left(n^{-1 / 2} T^{-1 / 2} h^{-1 / 2}\right) .
$$

Proof. Consider

$$
\begin{aligned}
& \frac{\sum_{t=1}^{T} f_{j t}\left[\widehat{\lambda}_{1 t}\left(j, x_{j}\right)-f_{u t}-\sum_{k \neq j} f_{k t} \widetilde{\lambda}_{2}\left(j, k, x_{j}\right)\right]}{\sum_{t=1}^{T} f_{j t}^{2}}-g_{j}\left(x_{j}\right) \\
= & \frac{1}{\sum_{t=1}^{T} f_{j t}^{2}} \sum_{t=1}^{T} f_{j t}\left\{\frac{\sum_{i=1}^{n} K_{h}\left(X_{j i}-x_{j}\right)\left[y_{i t}-f_{u t}-\sum_{k \neq j} f_{k t} \widetilde{g}_{k}\left(X_{k i}\right)\right]}{\sum_{i=1}^{n} K_{h}\left(X_{j i}-x_{j}\right)}-f_{j t} g_{j}\left(x_{j}\right)\right\} \\
= & \frac{1}{\sum_{t=1}^{T} f_{j t}^{2}} \sum_{t=1}^{T} f_{j t}\left\{\frac{\sum_{i=1}^{n} K_{h}\left(X_{j i}-x_{j}\right) \varepsilon_{i t}}{\sum_{i=1}^{n} K_{h}\left(X_{j i}-x_{j}\right)}+\frac{\sum_{i=1}^{n} K_{h}\left(X_{j i}-x_{j}\right)\left[g_{j}\left(X_{j i}\right)-g_{j}\left(x_{j}\right)\right]}{\sum_{i=1}^{n} K_{h}\left(X_{j i}-x_{j}\right)}\right\} \\
& -\frac{1}{\sum_{t=1}^{T} f_{j t}^{2}} \sum_{t=1}^{T} f_{j t} \frac{\sum_{i=1}^{n} K_{h}\left(X_{j i}-x_{j}\right) \sum_{k \neq j} f_{k t}\left[\widetilde{g}_{k}\left(X_{k i}\right)-g_{k}\left(X_{k i}\right)\right]}{\sum_{i=1}^{n} K_{h}\left(X_{j i}-x_{j}\right)} \\
= & U_{n 1}+U_{n 2}+U_{n 3} .
\end{aligned}
$$

The term $U_{n 2}$ is a standard bias term of order $h^{2}$. The term $U_{n 3}$ can be shown to be $O_{p}\left(h^{2}\right)+$ $o_{p}\left(n^{-1 / 2} T^{-1 / 2} h^{-1 / 2}\right)$ as in Linton (1997), where the $O_{p}\left(h^{2}\right)$ is a bias term. Interchanging summations and approximating $\sum_{i=1}^{n} K_{h}\left(X_{j i}-x_{j}\right) / n$ by $p_{j}\left(x_{j}\right)$ we obtain an approximation to the leading term $U_{n 1}$, i.e., $U_{n 1}=\widetilde{U}_{n 1}\left(1+O_{p}\left(n^{-1 / 2} h^{-1 / 2}\right)\right)$. The term $\widetilde{U}_{n 1}$ is a sum of independent random variables and is $O_{p}\left(n^{-1 / 2} h^{-1 / 2} T^{-1 / 2}\right)$ with mean zero and variance as stated in the theorem. Specifically, $E\left[\widetilde{\varepsilon}_{i} \mid X_{j i}=x_{j}\right]=0$ and

$$
\operatorname{var}\left[\widetilde{\varepsilon}_{i} \mid X_{j i}=x_{j}\right]=\frac{\sum_{t=1}^{T} f_{j t}^{2} \sigma_{j t}^{2}\left(x_{j}\right)}{\left(\sum_{t=1}^{T} f_{j t}^{2}\right)^{2}} \leq \frac{C}{T}
$$

for some $C<\infty$ for large enough $T$. Therefore, we can apply the Lindeberg CLT to $\widetilde{U}_{n 1}$. Furthermore, $\widetilde{U}_{n 1 j}\left(x_{j}\right)$ and $\widetilde{U}_{n 1 k}\left(x_{k}\right)$ are asymptotically independent by standard arguments for kernels with bounded support.

Now define $\widehat{f}_{t}^{[1]}$ as in (19). Using the above expansion it can be shown that the results of Lemma 2 and 4 continue to hold with $\widehat{f}_{t}^{[1]}$ replacing $\widetilde{f}_{t}$ with a difference sequence $\delta_{n, t}$. Then we can show that the conclusion of Lemma 5 and 6 continue to hold with $\widehat{g}_{j}^{[2]}$ replacing $\widehat{g}_{j}^{[1]}$. This process can be continued for any finite number of iterations, see Linton, Nielsen, and Van der Geer (2004). The only thing that changes in the expansions is the bias function, although the bias at each iteration can still be approximated by some bounded continuous function.

Proof of Theorem 2. By the triangle inequality

$$
\sup _{\theta \in \Theta}\left\|\widehat{M}_{T}(\theta)-M(\theta)\right\| \leq \sup _{\theta \in \Theta}\left\|\widehat{M}_{T}(\theta)-M_{T}(\theta)\right\|+\sup _{\theta \in \Theta}\left\|M_{T}(\theta)-M(\theta)\right\|,
$$

where $M_{T}(\theta)=T^{-1} \sum_{t=k+1}^{T} \psi\left(F_{t}, Z_{t} ; \theta\right)$. Applying uniform laws of large numbers for classes of smooth functions (Andrews (1987)) we obtain that $\sup _{\theta \in \Theta}\left\|M_{T}(\theta)-M(\theta)\right\|=o_{p}(1)$. Furthermore, by a first order expansion for each $\ell=1, \ldots, q$,

$$
\widehat{M}_{T \ell}(\theta)-M_{T \ell}(\theta)=\frac{1}{T} \sum_{t=k+1}^{T} \frac{\partial \psi_{\ell}}{\partial F_{t}}\left(\bar{F}_{t}, Z_{t} ; \theta\right)\left(\widehat{F}_{t}-F_{t}\right)
$$

where $\bar{F}_{t}$ are intermediate values. We next use the inequality $\operatorname{Pr}[C] \leq \operatorname{Pr}[C \cap D]+\operatorname{Pr}\left[D^{c}\right]$ with $D=\left\{\max _{1 \leq t \leq T}\left|\widehat{f}_{t}-f_{t}\right| \leq \delta_{n, T}\right\}$ and $\delta_{n, T} \rightarrow 0$ chosen such that $\operatorname{Pr}\left[D^{c}\right] \rightarrow 0$; this allows us to restrict attention to the event $D$. It follows that on this set by crude bounding using assumption B3 and (28) we have for some $C<\infty$,

$$
\begin{aligned}
& \left\|\frac{1}{T} \sum_{t=k+1}^{T} \frac{\partial \psi_{\ell}}{\partial F_{t}}\left(\bar{F}_{t}, Z_{t} ; \theta\right)\left(\widehat{F}_{t}-F_{t}\right)\right\| \\
\leq & C\left(\max _{1 \leq t \leq T}\left\|\widehat{f_{t}}-f_{t}\right\|\right) \frac{1}{T} \sum_{t=k+1}^{T}\left\|\theta-\theta_{0}\right\| \leq \delta,\left\|F^{\prime}-F_{t}\right\| \leq \delta_{n, T} \\
= & o_{p}(1) .
\end{aligned}
$$

It follows that $\sup _{\theta \in \Theta}\left\|\widehat{M}_{T}(\theta)-M_{T}(\theta)\right\|=o_{p}(1)$. Finally, the unique minimum condition B2 implies that $\widehat{\theta}-\theta_{0}=o_{p}(1)$.

By a Taylor expansion

$$
\begin{equation*}
\widehat{M}_{T}(\widehat{\theta})=\widehat{M}_{T}\left(\theta_{0}\right)+\frac{\partial \widehat{M}_{T}}{\partial \theta}(\bar{\theta})\left(\widehat{\theta}-\theta_{0}\right) \tag{45}
\end{equation*}
$$

where $\bar{\theta}$ are intermediate values. Furthermore, for each $\ell=1, \ldots, q$,

$$
\begin{aligned}
\widehat{M}_{T \ell}\left(\theta_{0}\right)= & M_{T \ell}\left(\theta_{0}\right)+\frac{1}{T} \sum_{t=k+1}^{T} \frac{\partial \psi_{\ell}}{\partial F_{t}}\left(F_{t}, Z_{t} ; \theta_{0}\right)\left(\widehat{F}_{t}-F_{t}\right) \\
& +\frac{1}{2 T} \sum_{t=k+1}^{T}\left(\widehat{F}_{t}-F_{t}\right)^{\top} \frac{\partial^{2} \psi_{\ell}}{\partial F_{t} \partial F_{t}^{\top}}\left(\bar{F}_{t}, Z_{t} ; \theta_{0}\right)\left(\widehat{F}_{t}-F_{t}\right),
\end{aligned}
$$

where $\bar{F}_{t}$ are intermediate values. By substituting in the expansion for $\widehat{F}_{t}-F_{t}$ it is easy to see that

$$
\begin{equation*}
\frac{1}{T} \sum_{t=k+1}^{T} \frac{\partial \psi_{\ell}}{\partial F_{t}}\left(F_{t}, Z_{t} ; \theta_{0}\right)\left(\widehat{F}_{t}-F_{t}\right)=O_{p}\left((n T)^{-1 / 2}\right)+O_{p}\left(h^{2}\right) \tag{46}
\end{equation*}
$$

We next use the inequality $\operatorname{Pr}[C] \leq \operatorname{Pr}[C \cap D]+\operatorname{Pr}\left[D^{c}\right]$ with $D=\left\{\max _{1 \leq t \leq T}\left|\widehat{f}_{t}-f_{t}\right| \leq \delta_{n, T} n^{-1 / 2}(\log T)^{\rho}\right\}$ and $\delta_{n, T} \rightarrow 0$ chosen such that $\operatorname{Pr}\left[D^{c}\right] \rightarrow 0$; this allows us to restrict attention to the event $D$. It
follows that on this set by crude bounding using assumption B3 we have for some $C<\infty$,

$$
\begin{aligned}
& \left\|\frac{1}{2 T} \sum_{t=k+1}^{T}\left(\widehat{F}_{t}-F_{t}\right)^{\top} \frac{\partial^{2} \psi_{\ell}}{\partial F_{t} \partial F_{t}^{\top}}\left(\bar{F}_{t}, Z_{t} ; \theta_{0}\right)\left(\widehat{F}_{t}-F_{t}\right)\right\| \\
\leq & C\left(\max _{1 \leq t \leq T}\left\|\widehat{f}_{t}-f_{t}\right\|\right)^{2} \frac{1}{T} \sum_{t=k+1}^{T} \sup _{\left\|\theta-\theta_{0}\right\| \leq \delta_{n, T},\left\|F^{\prime}-F_{t}\right\| \leq \delta_{n, T}}\left\|\frac{\partial^{2} \psi_{\ell}}{\partial F_{t} \partial F_{t}^{\top}}\left(F^{\prime}, Z_{t} ; \theta\right)\right\| \\
= & O_{p}\left(n^{-1}(\log T)^{2 \rho}\right)=o_{p}\left(T^{-1 / 2}\right) .
\end{aligned}
$$

Therefore, $\widehat{M}_{T}\left(\theta_{0}\right)=M_{T}\left(\theta_{0}\right)+o_{p}\left(T^{-1 / 2}\right)$. Similarly

$$
\frac{\partial \widehat{M}_{T}}{\partial \theta}(\bar{\theta})=\frac{\partial M_{T}}{\partial \theta}\left(\theta_{0}\right)+o_{p}(1)=\Gamma_{0}+o_{p}(1)
$$

In conclusion, we have $o_{p}\left(T^{-1 / 2}\right)=M_{T}\left(\theta_{0}\right)+\Gamma_{0}\left(\widehat{\theta}-\theta_{0}\right)$, and the result follows from arguments of Pakes and Pollard (1989, pp 1041-1042). In particular, a CLT for stationary mixing random variables is applied to $\sqrt{T} M_{T}\left(\theta_{0}\right)$, whence $\sqrt{T}\left(\widehat{\theta}-\theta_{0}\right)$ is asymptotically normal as stated.

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Table 1: Sample Statistics of Raw Security Characteristics

| Year | Firms | Mean |  |  |  | Standard Deviation |  |  |  | Skewness |  |  |  | Excess kurtosis |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | Size | Value | Mom | Vol | Size | Value | Mom | Vol | Size | Value | Mom | Vol | Size | Value | Mom | Vol |
| 1965 | 467 | 4.84 | 0.69 | 0.01 | 0.06 | 1.69 | 0.40 | 0.02 | 0.03 | -0.25 | 1.04 | 0.18 | 1.21 | -0.14 | 0.47 | 0.59 | 1.04 |
| 1970 | 1562 | 4.62 | 0.52 | 0.00 | 0.10 | 1.56 | 0.29 | 0.03 | 0.04 | 0.23 | 0.82 | 0.23 | 0.74 | -0.47 | 0.19 | 0.46 | 0.09 |
| 1975 | 3394 | 3.44 | 1.40 | -0.01 | 0.14 | 1.81 | 0.90 | 0.03 | 0.06 | 0.47 | 0.91 | 0.03 | 0.59 | -0.38 | 0.27 | 0.19 | -0.14 |
| 1980 | 3465 | 3.72 | 1.21 | 0.02 | 0.12 | 1.86 | 0.69 | 0.03 | 0.06 | 0.20 | 0.79 | 0.44 | 0.69 | -0.38 | 0.15 | 0.40 | 0.00 |
| 1985 | 4017 | 4.20 | 0.73 | -0.02 | 0.10 | 1.95 | 0.43 | 0.03 | 0.05 | 0.17 | 0.69 | -0.37 | 0.76 | -0.42 | -0.04 | 0.22 | 0.06 |
| 1990 | 4595 | 4.13 | 0.76 | 0.01 | 0.11 | 2.18 | 0.48 | 0.03 | 0.06 | 0.24 | 0.91 | -0.02 | 1.00 | -0.33 | 0.47 | 0.59 | 0.39 |
| 1995 | 5583 | 4.78 | 0.62 | 0.01 | 0.11 | 2.00 | 0.40 | 0.03 | 0.06 | 0.37 | 0.91 | -0.07 | 0.94 | -0.31 | 0.28 | 0.57 | 0.33 |
| 2000 | 6003 | 5.25 | 0.65 | 0.01 | 0.17 | 2.12 | 0.46 | 0.04 | 0.09 | 0.38 | 1.01 | 0.55 | 0.90 | -0.27 | 0.47 | 0.73 | 0.31 |
| 2005 | 4952 | 6.06 | 0.52 | 0.03 | 0.11 | 2.03 | 0.31 | 0.03 | 0.07 | 0.23 | 0.85 | 0.58 | 1.08 | -0.28 | 0.39 | 0.66 | 0.42 |
| Med | 4017 | 4.62 | 0.69 | 0.01 | 0.11 | 1.95 | 0.43 | 0.03 | 0.06 | 0.23 | 0.91 | 0.18 | 0.90 | -0.33 | 0.28 | 0.57 | 0.31 |

Table 1 shows some descriptive statistics of the cross-sectional data for July at five-year intervals: the number of securities in the annual cross section, and the first four crosssectional moments of the four raw characteristics. Separately provided are the time series medians over the full 42 year period, also using July data.

Figure 1: Time Series Plots of Cross-sectional Correlations between the Characteristics


Figure 1 illustrates the time varying nature of the correlations between characteristics by showing the cross-sectional correlation for each pair of characteristics each July.

## Table 2: Characteristic-Beta Functions and Standard Errors



Table 2 shows the estimates of the characteristic-beta functions at some selected characteristic values and the heteroskedasticity-consistent standard errors for each of these estimates. The table reports results for each characteristic over a support ranging from the empirical $2.5 \%$ to the $97.5 \%$ quantile.

Figure 2: The Characteristic-beta Functions

## Size characteristic-beta function



Momentum characteristic-beta function


Value characteristic-beta function


Own-Volatility characteristic-beta function


The solid lines in Figure 2 display the estimated characteristic-beta functions over all grid points. For comparative purposes the dashed line depicts the linear Rosenberg-type model with the same identification conditions imposed.

Figure 3: Characteristic-Beta Functions on Four 126-Month Subperiods

Size characteristic-beta functions



Momentum characteristic-beta functions


Value characteristic-beta functions


Own-vol characteristic-beta functions


Figure 3 shows the characteristic-beta functions from four subsamples. SP1 line is a function for 1963 July-1973 December period, SP2 for 1974 January-1984 June period, while SP3 and SP4 for 1984 July1994 December and 1995 January-2005 July periods, respectively.

Table 3: Uncentered R-squared Statistics Using Subsets of the Characteristics

|  | Marginal UR2 Statistics when Adding Factors First or Last to the Model |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
|  | Market | Size | Value | Momentum | Volatility |
| Adding first | $12.23 \%$ | 1.50\% | 0.66\% | 1.26\% | 2.01\% |
| Adding last | 12.12\% | 0.89\% | 0.34\% | 0.74\% | 1.09\% |
|  | UR2 Linear and Nonlinear Five Factor Model |  |  |  |  |
|  | Linear | Nonlinear |  |  |  |
| UR2 | 16.20\% | 16.44\% |  |  |  |

Table 3 shows the time-series averages of uncentered cross-sectional $R^{2}$ (UR2) statistics as a measure of the explanatory power of the factor model. The upper part of the table shows average UR2 statistics from cross-sectional regressions of excess returns on each characteristic-beta function singly as well as their marginal contribution given the other four. The lower part of the table shows the UR2 statistic based on the model with all five characteristic-beta functions.

Table 4: Factor Return Statistics and Comparison to Fama French Factor Mimicking Portfolios

|  | Panel a: Factor Return Statiatica |  |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | Market | Size | Value | Momentum | Volatility |  |  |  |  |
| Annualized Mean | 9.65\% | $-2.66 \%$ | 0.80\% | -0.28\% | -1.27\% |  |  |  |  |
| Annualized Volatility | 19.47\% | $5.78 \%$ | 3.54\% | 5.85\% | 7.29\% |  |  |  |  |
| \% Periods significant* | 91.07\% | 74.40\% | 57.34\% | $66.47 \%$ | $75.20 \%$ |  |  |  |  |
| overall p-value | 0.00\% | 0.00\% | 0.00\% | 0.00\% | 0.00\% |  |  |  |  |
|  |  | Pancl b | irical Fa | Return and | ama French | micking P | Retur | lations |  |
|  | Market | Size | Value | Momentum | Volatility | RMRF | SMB | HML | FF_MOM |
| Market | 1.00 | 0.00 | -0.16 | -0.39 | 0.80 | 0.86 | 0.66 | -0.30 | -0.17 |
| Size | 0.00 | 1.00 | 0.19 | 0.15 | -0.01 | 0.34 | -0.34 | -0.10 | -0.11 |
| Value | -0.16 | 0.19 | 1.00 | -0.12 | -0.32 | -0.23 | -0.18 | 0.78 | -0.21 |
| Momentum | -0.39 | 0.15 | -0.12 | 1.00 | -0.41 | -0.17 | -0.18 | -0.08 | 0.80 |
| Volatility | 0.80 | -0.01 | -0.32 | -0.41 | 1.00 | 0.65 | 0.63 | -0.44 | -0.17 |
| RMRF | 0.86 | 0.34 | -0.23 | -0.17 | 0.65 | 1.00 | 0.30 | -0.41 | -0.06 |
| SMB | 0.66 | -0.34 | -0.18 | -0.18 | 0.63 | 0.30 | 1.00 | -0.28 | 0.02 |
| HML | -0.30 | -0.10 | 0.78 | -0.08 | -0.44 | -0.41 | -0.28 | 1.00 | -0.12 |
| FF_MOM | -0.17 | -0.11 | -0.21 | 0.80 | -0.17 | -0.06 | 0.02 | -0.12 | 1.00 |

[^9]Figure 4: Characteristic-based Mispricing Function


Own-Vol Characteristic Mispricing


Figure 4 shows the four additive nonparametric characteristic-based mispricing functions over the grid of points between -3 and 3 at equally-spaced intervals of 0.1 .

Table 5: Vector Autoregression for the Factor Returns

|  | Market | Size | Value | Momentum | Own- <br> Volatility |
| :---: | :---: | :---: | :---: | :---: | :---: |
|  |  |  |  |  |  |
| Market(-1) | 0.122 | -0.074 | -0.002 | -0.014 | -0.005 |
|  | $(0.086)$ | $(0.002)$ | $(0.001)$ | $(0.016)$ | $(0.020)$ |
| Size(-1) | 0.733 | 0.133 | -0.001 | -0.116 | 0.171 |
|  | $(0.239)$ | $(0.047)$ | $(0.004)$ | $(0.046)$ | $(0.056)$ |
| Value(-1) | -2.180 | -0.496 | 0.292 | -0.475 | -1.220 |
|  | $(3.800)$ | $(0.751)$ | $(0.057)$ | $(0.723)$ | $(0.889)$ |
| Mom(-1) | 0.201 | -0.144 | -0.005 | -0.073 | 0.118 |
|  | $(0.287)$ | $(0.057)$ | $(0.004)$ | $(0.055)$ | $(0.067)$ |
| Own-Vol (-1) | 0.244 | -0.047 | 0.008 | -0.123 | 0.200 |
|  | $(0.431)$ | $(0.085)$ | $(0.006)$ | $(0.082)$ | $(0.101)$ |
| Constant | 0.008 | -0.001 | 0.000 | 0.000 | 0.000 |
|  | $(0.003)$ | $(0.001)$ | $(0.000)$ | $(0.001)$ | $(0.001)$ |
| R2 | 0.054 | 0.129 | 0.068 | 0.054 | 0.061 |

Table 5 shows the results from a first-order vector autoregression of the five factors returns on their lagged values. Heteroskedasticity-consistent standard errors are shown in parentheses below each estimated coefficient.

Table 6: Vector Autoregression for the Squared Factor Returns

|  | Market^2 | Size ${ }^{\wedge} 2$ | Value^2 | Momentum ${ }^{\text {^2 }}$ | Own-Volatility ${ }^{\wedge}$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| Market(-1)^2 | -0.002 | 0.004 | 0.000 | 0.001 | -0.004 |
|  | (0.052) | (0.002) | (0.000) | (0.003) | (0.004) |
| Size(-1) ${ }^{\wedge} 2$ | 1.400 | 0.238 | 0.002 | 0.280 | 0.520 |
|  | (1.370) | (0.052) | (0.000) | (0.070) | (0.090) |
| Value(-1) ${ }^{\wedge} 2$ | -76.030 | 32.390 | 0.230 | 14.080 | 19.070 |
|  | (248.300) | (9.480) | (0.062) | (12.610) | (16.370) |
| $\operatorname{Mom}(-1){ }^{\wedge} 2$ | 0.699 | -0.102 | 0.000 | 0.126 | 0.014 |
|  | (1.420) | (0.054) | (0.000) | (0.072) | (0.094) |
| Own-Vol (-1) ${ }^{\wedge} 2$ | -0.390 | -0.013 | 0.000 | -0.110 | 0.041 |
|  | (1.360) | (0.052) | (0.000) | (0.069) | (0.090) |
| Constant | 0.004 | 0.000 | 0.000 | 0.000 | 0.000 |
|  | (0.001) | (0.000) | (0.000) | (0.000) | (0.000) |
| R ${ }^{2}$ | 0.003 | 0.145 | 0.257 | 0.067 | 0.141 |

Table 6 shows a first-order vector autoregression of squared factor returns on their lagged values.
Heteroskedasticity-consistent standard errors are shown in parentheses below each estimated coefficient.

Figure 5: Time series of cross-sectional root-mean-square asset-specific return


Figure 5 shows the time-series of cross-sectional mean-square asset-specific return over the sample period July 1964 to June 2005.

Figure 6: The Characteristic-beta Functions with Observed Market Index


Size Characteristic-Beta Function

Momentum Characteristic-beta Function


Value Characteristic-beta Function


Own-Volatility Characteristic-beta Function


Figure 6 displays the estimated characteristic-beta functions over the range of grid points covering $95 \%$ of characteristic values. The model is estimated with an observed market index. The characteristic-beta functions are standardized to have a capitalization-weighted mean of zero and capitalization-weighted variance of one.

Table 7: Tests of the Capital Asset Pricing Model as a Restriction on the Non-Market Factor Returns

|  | intercept |  | beta |
| :---: | :---: | :---: | :---: |
|  | Full Period: $7 / 1963-6 / 2005$ |  |  |
|  |  | R2 |  |
| Size | $-.0011(-3.40)$ | $.0065(.875)$ | 0.002 |
| Value | $.0006(2.09)$ | $-.057(-9.06)$ | 0.14 |
| Momentum | $-.0004(.755)$ | $-.039(-3.41)$ | 0.023 |
| Own-Volatility | $-.0008(-3.21)$ | $.105(17.0)$ | 0.367 |
| GRS (p-value) | $10.49(.0000)$ | First Subperiod: $7 / 1963-12 / 1973$ |  |

First Subperiod: 7/1963-12/1973

|  |  |  |  |
| :---: | :---: | :---: | :---: |
| Size | $-.0006(-1.07)$ | $-.061(-4.08)$ | 0.118 |
| Value | $.0001(.275)$ | $-.040(-2.92)$ | 0.065 |
| Momentum | $.0012(1.54)$ | $-.056(-2.68)$ | 0.055 |
| Own-Volatility | $-.0006(1.69)$ | $.104(10.4)$ | 0.465 |
| GRS (p-value) | $3.00(.021)$ |  |  |
|  | Second Subperiod: $1 / 1974-6 / 1984$ |  |  |


| Size | $-.0027(-3.78)$ | $-.010(-.723)$ | 0.004 |
| :---: | :---: | :---: | :---: |
| Value | $.0003(.499)$ | $-.063(-5.73)$ | 0.207 |
| Momentum | $-.0007(-.734)$ | $-.021(-1.09)$ | 0.009 |
| Own-Volatility | $-.0010(-2.97)$ | $.074(10.5)$ | 0.469 |
| GRS (p-value) | $11.57(.0000)$ |  |  |
|  | Third Subperiod: $7 / 1984-12 / 1994$ |  |  |

Third Subperiod: 7/1984-12/1994


Table 7 shows the results from univariate time-series regression of each of the four characteristic-based factors on a constant and the excess return to the market index. The CAPM implies that the set of four intercepts from these univariate regressions should jointly equal zero. T-statistics are shown in parentheses next to each estimated coefficient. The regressions are run over the full 504 month sample period and for each of four 126-month subperiods. The Gibbons-Ross-Shanken test statistic of the joint hypothesis that all four intercepts are zero is also shown, together with its p-value in parentheses.


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[^1]:    ${ }^{1}$ Bai assumes that the loadings $\beta_{i j}$ are fixed in repeated samples but treats $f_{j t}$ as random, although he remarks about the reverse case.

[^2]:    ${ }^{2}$ It is possible to include additional nonparametric functions based on other observed variates strictly exogenous relative to $y_{i t}$; this does not require any additional identification conditions beyond the mean-zero condition.

[^3]:    ${ }^{3}$ In our case, the cross-sectionally averaged returns are not informative except about $f_{u t}$, since $\bar{y}_{t}=f_{u t}+$ $\sum_{j=1}^{J} \bar{g}_{j} f_{j t}+\bar{\varepsilon}_{t}=f_{u t}+O_{p}\left(n^{-1 / 2}\right)$ under the assumption that $E\left[g_{j}\left(X_{j i}\right)\right]=0$, where $\bar{y}_{t}=\sum_{i=1}^{n} y_{i t} / n, \bar{\varepsilon}_{t}=\sum_{i=1}^{n} \varepsilon_{i t} / n$, and $\bar{g}_{j}=\sum_{i=1}^{n} g_{j}\left(X_{j i}\right) / n$.

[^4]:    ${ }^{4}$ The setting is formally that of a triangular array because the distribution of $\bar{\varepsilon}_{i}$ changes with $T(n)$.
    ${ }^{5}$ For identification we require that $\bar{f}_{j} \neq 0$ for all $j$, which is like in Pesaran (2006).

[^5]:    ${ }^{6}$ The rate is partly determined by our assumption that $\sum_{t=1}^{T} f_{j t}^{2} / T$ stays bounded away from zero and infinity, as is appropriate for a stationary time series. A more general theory can be written for the case where $\sum_{t=1}^{T} f_{j t}^{2}=O\left(T^{\alpha}\right)$ for any $\alpha \geq 0$. In our application the case $\alpha=1$ seems the most relevant.

[^6]:    ${ }^{7}$ See Bester and Hansen (2007) for a treatment of an even more general model allowing for nonparametric covariate effects and correlated unobserved random effects.
    ${ }^{8}$ Haag (2007) has analyzed a testing problem involving backfitting estimators. A number of different nonparametric test statistics based on standard kernel estimators are reviewed in Li and Racine (2007).

[^7]:    ${ }^{9}$ http://mba.tuck.dartmouth.edu/pages/faculty/ken.french/Data_Library/det_mom_factor.html

[^8]:    ${ }^{10}$ This model could arise from the general statistical factor model with relevant covariates, for example, suppose that $\beta_{j i}=g_{j}\left(X_{j i}\right)+\eta_{j i}$, where $\eta_{j i}$ are unobserved stochastic errors satisfying mean independence from the observed characteristics. The resulting model for returns is of the form (34).

[^9]:    The upper part of Table 4 shows the mean, volatility, and statistical significance of sach factor. The statistical significance is calculated as the percentage of aignificant $t$ atatiatics for each factor by estimating for each crosa-sectional regression, the t-atatiatic for each eatimated coefficient, based on Hansen-White heteroskedasticity consiaten otandard errore. Then for each factor, finding the sverage number of crose-sectional regresaion t-atatiotics that are significant at a $95 \%$ confidence level acrose the 504 time periods. The asgregate p-value is aloo provided. The lower part of Table 4 disploys the correlations between the eatimated factiona, olong with the three Fomo-French factors, RMRF, SMB, ond HML and a momentum factor FF-MOM croated by Ken French. RMRF is the Femo-French market factor, it is the return to the volue-weighted market index minus the riakfree return; SMB is the return to a amoll copitalization portfolio minus the return to a large-copitalization portfolio; HML is the return to a high book-toprice portfolio minus the return to o low book-to-price portfolio. FF-MOM is the return to a portfolio with high cumulative returns over the past twelve months minus the return to a portfolio with low cumulative returns over the past twelve monthe, adjusted to have roughly equal average capitalization.

