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## To cite this version:

Oliver B. Linton, Enno Mammen. Nonparametric transformation to white noise. Econometrics, 2007, 142 (1), pp.241. 10.1016/j.jeconom.2007.05.018 . hal-00501798

## HAL Id: hal-00501798

## https://hal.science/hal-00501798

Submitted on 12 Jul 2010

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## Author's Accepted Manuscript

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PII:
S0304-4076(07)00129-7
DOI:
doi:10.1016/j.jeconom.2007.05.018
Reference:
ECONOM 2958

To appear in: Journal of Econometrics

www.elsevier.com/locate/jeconom

Received date: 9 August 2006
Revised date: 2 April 2007
Accepted date: 3 May 2007
Cite this article as: Oliver B. Linton and Enno Mammen, Nonparametric transformation to white noise, Journal of Econometrics (2007), doi:10.1016/j.jeconom.2007.05.018

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# Nonparametric Transformation to White Noise 

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April 2, 2007


#### Abstract

We consider a semiparametric distributed lag model in which the "news impact curve" $m$ is nonparametric but the response is dynamic through some linear filters. A special case of this is a nonparametric regression with serially correlated errors. We propose an estimator of the news impact curve based on a dynamic transformation that produces white noise errors. This yields an estimating equation for $m$ that is a type two linear integral equation. We investigate both the stationary case and the case where the error has a unit root. In the stationary case we establish the pointwise asymptotic normality. In the special case of a nonparametric regression subject to time series errors our estimator achieves efficiency improvements over the usual estimators, see Xiao, Linton, Carroll, and Mammen (2003). In the unit root case our procedure is consistent and asymptotically normal unlike the standard regression smoother. We also present the distribution theory for the parameter estimates, which is non-standard in the unit root case. We also investigate its finite sample performance through simulation experiments.


Key words: Efficiency; Inverse Problem; Kernel Estimation; Nonparametric regression; Time Series; Unit Roots.

Journal of Economic Literature Classification: C14

[^0]
## 1 Introduction

In this paper we discuss the estimation of the unknown quantities in the model

$$
\begin{equation*}
B(L) Y_{t}=A(L) m\left(X_{t}\right)+\varepsilon_{t}, \tag{1}
\end{equation*}
$$

where $\varepsilon_{t}$ is a martingale difference sequence with respect to the past of $Y_{t}$ and current and past regressors $X_{t}$, while $A(L)=\sum_{j=0}^{\infty} a_{j} L^{j}$ and $B(L)=\sum_{j=0}^{\infty} b_{j} L^{j}$ are lag polynomial operators with $a_{0}=b_{0}=1$ for identification, where $L x_{t}=x_{t-1}$. The function $m($.$) is assumed to be unknown but$ smooth, and is the object of central interest, although the dynamics of the model represented by $A(L), B(L)$ are also fundamental to the interpretation.

We first discuss a special case of central interest, the nonparametric regression model

$$
\begin{equation*}
Y_{t}=m\left(X_{t}\right)+u_{t}, t=1, \ldots, T \tag{2}
\end{equation*}
$$

where the covariates follow some stationary mixing process, while the residual process $u_{t}$ satisfies

$$
\begin{equation*}
A(L) u_{t}=\varepsilon_{t}=\sum_{j=0}^{\infty} a_{j} u_{t-j} . \tag{3}
\end{equation*}
$$

In this case, $A(L) Y_{t}=A(L) m\left(X_{t}\right)+\varepsilon_{t}$, which is a special case of (1) with $A(L)=B(L)$. The parametric version of the regression model (2) and (3) is a standard teaching topic in graduate econometrics, Harvey (1981, Chapter 6). In the semiparametric model there are many standard estimators of $m$ and of the parameters of $A(L)$ that are consistent under summability conditions on $A$, see for example Robinson (1983), Bierens (1983), Masry and Fan (1997), Hidalgo (1997), and Fan and Yao (2003). However, unlike in the parametric case, the standard kernel regression smoothers do not take account of the correlation structure in $X_{t}$ or $u_{t}$ and estimate the regression function in the same way as if these processes were independent. Furthermore, the variance of such estimators is proportional to the short run variance of $u_{t}, \sigma_{u}^{2}=\operatorname{var}\left(u_{t}\right)$ and does not depend on the regressor or error covariance functions $\operatorname{cov}\left(X_{t}, X_{t-j}\right), \operatorname{cov}\left(u_{t}, u_{t-j}\right), j \neq 0$. This is a bit surprising in comparison with the parametric case. One might think that there is useful information in the autocorrelation structure for estimation of the mean. This point has been addressed recently by Xiao, Linton, Carroll, and Mammen (2003) who proposed a more efficient estimator of $m$ based on a prewhitening transformation

$$
\begin{equation*}
Y_{t}-\sum_{j=1}^{\infty} a_{j}\left(Y_{t-j}-m\left(X_{t-j}\right)\right)=m\left(X_{t}\right)+\varepsilon_{t} \tag{4}
\end{equation*}
$$

where the right hand side is now a standard nonparametric regression with whitened errors. The transform implicitly takes account of the autocorrelation structure. In practice they replaced the
unknown quantities on the left hand side by preliminary estimates of $m$ and $a_{j}(\alpha)$. Their procedure improves in terms of variance over the usual kernel smoothers.

The model (1) is more general than nonparametric regression with autocorrelated errors and is perhaps more rightly viewed as a generalization of the distributed lag model. The traditional distributed lag model (with $m(x)=x$ ) has been very popular in economics, Dhrymes (1971). ${ }^{1}$ More recently, Hendry, Pagan, and Sargan (1984) reviewed the specification of such models and gave a taxonomy of special cases. It can be motivated from some simple economic relationships being distorted by adaptive expectations, partial adjustment, etc., see Harvey (1981, Chapter 7). Suppose there is a latent variable $Y^{*}$ that has some equilibrium relationship with covariate $X$, which in general can be nonlinear so that $Y_{t}^{*}=m\left(X_{t}\right)$. Then suppose that actual $Y$ only responds to $Y^{*}$ with some lagging mechanism, for example, $Y_{t}-Y_{t-1}=\gamma\left[Y_{t}^{*}-Y_{t-1}\right]+\varepsilon_{t}$ for some $\gamma \in(0,1)$, then we obtain a special case of (1). ${ }^{2}$ The lags arise because production takes time or because agents take to respond to a signal or because there are institutional constraints. The traditional applications were in for example production studies where $Y_{t}$ is output and $X_{t}$ is the capital/labour ratio of a given firm or industry observed over time. More recent applications have been in rational expectations models where the data are at different frequencies, Hansen and Hodrick (1980). The issues concerning formulation and estimation of the lag polynomials $A, B$ are pretty much resolved in the linear case, see Hannan and Deistler (1988) for a more recent discussion in the multivariate case. Linearity of $m$ is just a convenience and was adopted many years ago when computational and technical issues were binding. We allow for nonlinear $m$ because for some problems linear $m$ is not well motivated and at odds with the data. Note that model (1) includes as a special case the so-called NARMAX model introduced in Chen and Billings (1989) and used frequently by systems engineers in which the function $m$ is approximated by some polynomial with unknown coefficients. Finally, we remark that the $\operatorname{ARCH}(\infty)$ model of Linton and Mammen (2005) is a special case when $Y_{t}=y_{t}^{2}$ and $X_{t}=y_{t-1}$, while $B(L)=1$. This model has been treated elsewhere.

We treat only the case where $A(L), B(L)$ are described by a finite dimensional parameter $\theta=$ $(\alpha, \beta) \in \mathbb{R}^{p}$ with $\alpha \in \mathbb{R}^{p_{a}}$ parameterizing $A$ and $\beta \in \mathbb{R}^{p_{b}}$ parameterizing $B$. We propose a strategy for estimation of $m$ along with the parameters of $A(L)$ in (2), (3). This is essentially to estimate the

[^1]transformed model (1) as an additive (possibly infinite order) nonparametric regression, see Hastie and Tibshirani (1991). Recently, Linton and Mammen (2005) have shown how to estimate similar models using the theory of linear integral equations of the second kind; see also Carrasco, Florens and Renault (2006). We obtain an estimating equation for $m$ that is a type two linear integral equation for each parameter value $\theta$. To obtain the parameters $\theta$ we optimize a profile likelihood criterion. We show that our method has attractive theoretical and finite sample properties. In particular, in the special case of nonparametric regression with autocorrelated error it has smaller asymptotic variance than the main method of Xiao, Linton, Carroll, and Mammen (2003). Furthermore, the asymptotics require weaker conditions with regard to the memory properties of the error terms. We define our method in the general model (1). In that case there is not such an obvious alternative estimator of the function $m$. We mostly consider the case where both $X_{t}, Y_{t}$ are stationary and mixing processes in which case the main statistical issue is efficiency. We also consider the case where some of the variables are nonstationary. This could arise for example from a unit root in the residual $u_{t}$ or in $X_{t}$ or in both, see Phillips and Park (1998). In this case, estimating in the original data (2) may lead to inconsistency, whereas the transformation involved in (1) yields error terms with a lower order of nonstationarity/persistence and hence consistency can be obtained, see Marinucci and Robinson (2003). The estimation method is more or less the same as in the stationary case although the justification of it differs. The distribution theory for the parametric part though is non standard in this case: in fact we obtain $T$ convergence to the Dickey-Fuller distribution under the unit root.

## 2 The Stationary Case

In this section we suppose that $\left(Y_{t}, X_{t}\right)$ are jointly stationary and weakly dependent mixing processes and describe our estimation methods and their properties for this case.

### 2.1 Estimation Method

### 2.1.1 Population Characterization

We first suppose that $A(L), B(L)$ are known. Letting $Z_{t}=B(L) Y_{t}$ we have

$$
Z_{t}=A(L) m\left(X_{t}\right)+\varepsilon_{t}=\sum_{j=0}^{\infty} a_{j} m\left(X_{t-j}\right)+\varepsilon_{t}
$$

which is an additive autoregression with i.i.d. errors where the additive components are subject to the restriction that they all share a common function $m$. In view of the assumed stationarity, define
the function $m$ as the minimizer of the criterion

$$
\begin{equation*}
Q\left(\theta_{0}, m\right)=E\left[\left\{Z_{0}-\sum_{j=0}^{\infty} a_{j} m\left(X_{-j}\right)\right\}^{2}\right] \tag{5}
\end{equation*}
$$

This problem can be viewed as a projection in a suitable Hilbert space. Let $L_{2}\left(f_{0}\right)$ be the Hilbert space of square integrable functions with respect to the marginal density of $X_{t}$ denoted $f_{0}$. For the stationary mixing process $\left\{X_{t}\right\}_{t=-\infty}^{\infty}$, provided $\sum_{j=0}^{\infty}\left|a_{j}\right|<\infty$, the random variable $\sum_{j=0}^{\infty} a_{j} m\left(X_{-j}\right)$ is square integrable for any function $m \in L_{2}\left(f_{0}\right)$. The set $\mathcal{G}=\left\{\sum_{j=0}^{\infty} a_{j} m\left(X_{-j}\right): m \in L_{2}\left(f_{0}\right)\right\}$ can be viewed as a subspace of the Hilbert space of square integrable functions defined on the infinite product of random variables $\underline{X}=\left(X_{0}, X_{-1}, \ldots\right)$. By the projection theorem there exists a unique member of $\mathcal{G}$ closest to the random variable $Z_{0}$.

A necessary condition for $m$ to be the minimizer of (5) is that it satisfies the first order condition

$$
\begin{equation*}
E\left[\left\{Z_{0}-\sum_{j=0}^{\infty} a_{j} m\left(X_{-j}\right)\right\} \sum_{k=0}^{\infty} a_{k} h\left(X_{-k}\right)\right]=0 \tag{6}
\end{equation*}
$$

for any measurable function $h$ for which this expectation is well-defined. See Sagan (1969), Theorem 1.7 for example. The second order condition is $-E\left[\left\{\sum_{k=0}^{\infty} a_{k} h\left(X_{-k}\right)\right\}^{2}\right]$ which is negative implying that the solution of the first order condition does indeed (locally) minimize the criterion. Taking $h($. to be the Dirac delta function, we have that

$$
\begin{equation*}
\sum_{j=0}^{\infty} a_{j} E\left[Z_{0} \mid X_{-j}=x\right]=\sum_{j=0}^{\infty} a_{j}^{2} m(x)+\sum_{j \neq k} \sum_{j} a_{j} a_{k} E\left[m\left(X_{-j}\right) \mid X_{-k}=x\right] \tag{7}
\end{equation*}
$$

for each $x .^{3}$ This is an implicit equation for $m($.$) . It can be re-expressed as a linear type two integral$ equation in $L_{2}\left(f_{0}\right)$. Define $a_{j}^{\dagger}=a_{j} / \sum_{l=0}^{\infty} a_{l}^{2}$ and $a_{j}^{*}=\sum_{k \neq 0} a_{j+k} a_{k} / \sum_{l=0}^{\infty} a_{l}^{2}$, and let $f_{0, j}$ be the joint density of ( $X_{t}, X_{t-j}$ ). Then

$$
\begin{gather*}
m(x)=m^{*}(x)+\int \mathcal{H}(x, y) m(y) f_{0}(y) d y, \text { or } m=m^{*}+\mathcal{H} m  \tag{8}\\
m^{*}(x)=\sum_{j=0}^{\infty} a_{j}^{\dagger} E\left[Z_{0} \mid X_{-j}=x\right] \\
\mathcal{H}(x, y)=-\sum_{j= \pm 1}^{ \pm \infty} a_{j}^{*} \frac{f_{0, j}(y, x)}{f_{0}(y) f_{0}(x)}
\end{gather*}
$$

[^2]This is similar to the equation derived in Linton and Mammen (2005) with the exception that there $X_{t}$ was lagged values of $Y_{t}$. Equation (8) is an implicit equation in $m$ and we need some conditions on the operator $\mathcal{H}(x, y)$ to guarantee that there exists a unique solution.

Assumption A1. The operator $\mathcal{H}(x, y)$ satisfies the Hilbert-Schmidt condition i.e.,

$$
\iint \mathcal{H}(x, y)^{2} f_{0}(x) f_{0}(y) d x d y<\infty
$$

A sufficient condition for A1 is that the joint densities $f_{0, j}(y, x)$ have compact support and $f_{0}(x)$ is bounded away from zero on this support, which we shall assume below. However, this is not necessary and condition A1 can hold for many covariate processes with unbounded support. We shall however restrict attention to the case where the support of the marginal covariate density $f_{0}$ is a compact set $[\underline{x}, \bar{x}]$. Then the operator $\mathcal{H}$ is a bounded compact linear operator on the Hilbert space of functions $L_{2}\left(f_{0}\right)$. It is also self-adjoint, see Linton and Mammen (2005). It therefore has a countable number of eigenvalues ${ }^{4}$ :

$$
\infty>\left|\lambda_{1}\right| \geq\left|\lambda_{2}\right| \geq \ldots
$$

with $\sum_{j=0}^{\infty} \lambda_{j}^{2}<\infty$. Also, the value 0 is a cluster point of the set $\left\{\lambda_{j}\right\}_{j=1}^{\infty}$ and 0 is the only cluster point, see Kress (1999, Theorem 3.9). The spectral radius of an operator $\mathcal{H}$ is defined as $r(\mathcal{H})=\sup _{j}\left|\lambda_{j}\right|$, and in this case it is finite.

Assumption A2. There exists no measurable function $m($.$) with \int m(x)^{2} f_{0}(x) d x=1$ such that $\sum_{j=0}^{\infty} a_{j} m\left(X_{t-j}\right)=0$ with probability one.

This condition rules out a certain 'concurvity' in the stochastic process $\left\{X_{t}\right\}$. That is, the data cannot be functionally related in this particular way. In the $\operatorname{AR}(1)$ case this says that there are no nontrivial functions $m$ that satisfy $m\left(X_{t}\right)-\rho m\left(X_{t-1}\right)=0$ with probability one. ${ }^{5}$ A consequence of A2 is that $\sup _{j} \lambda_{j}<1$ and therefore the operator $I-\mathcal{H}$ is strictly positive definite. Therefore, there exists a unique solution to (8) that satisfies

$$
\begin{equation*}
m=(I-\mathcal{H})^{-1} m^{*} \tag{9}
\end{equation*}
$$

This is the main characterization used for estimation, although we must first extend this to the case where a general $\theta$ is used not necessarily the true $\theta_{0}$.

For each $\theta=(\alpha, \beta) \in \Theta$, define $Z_{t}(\beta)=\sum_{j=0}^{\infty} b_{j}(\beta) Y_{t-j}$ and $g_{j}(x ; \beta)=E\left[Z_{t}(\beta) \mid X_{t-j}=x\right]$, $j=0, \pm 1, \ldots$

$$
m_{\theta}^{*}(x)=\sum_{j=0}^{\infty} a_{j}^{\dagger}(\alpha) g_{j}(x ; \beta)
$$

[^3]\[

$$
\begin{equation*}
\mathcal{H}_{\theta}(x, y)=-\sum_{j= \pm 1}^{ \pm \infty} a_{j}^{*}(\alpha) \frac{f_{0, j}(y, x)}{f_{0}(y) f_{0}(x)} \tag{10}
\end{equation*}
$$

\]

where $a_{j}^{\dagger}(\alpha)=a_{j}(\alpha) / \sum_{l=0}^{\infty} a_{l}^{2}(\alpha)$ and $a_{j}^{*}(\alpha)=\sum_{k \neq 0} a_{j+k}(\alpha) a_{k}(\alpha) / \sum_{l=0}^{\infty} a_{l}^{2}(\alpha)$. We now let $m$ vary with $\theta$, that is, (5) is defined for any $\theta$, and let $m_{\theta}$ be the function that minimizes (5); this satisfies $m_{\theta}=\left(I-\mathcal{H}_{\theta}\right)^{-1} m_{\theta}^{*}$ for all $\theta$ provided the conditions A1 and A2 hold uniformly over the parameter space $\Theta$. Furthermore, we can define $\theta=\theta_{0}$ as the minimizer of

$$
\begin{equation*}
Q\left(\theta, m_{\theta}\right)=E\left[\left\{Z_{0}(\beta)-\sum_{j=0}^{\infty} a_{j}(\alpha) m_{\theta}\left(X_{-j}\right)\right\}^{2}\right] \tag{11}
\end{equation*}
$$

with respect to $\theta \in \Theta$. Let $m_{0}=m_{\theta_{0}}$. We adopt this profiling approach to defining $\theta_{0}, m_{0}$ as this is the way our estimation strategy works. We suppose that assumptions A1 and A2 hold uniformly over the parameter space $\Theta$ so that for each $\theta \in \Theta, m_{\theta}=\left(I-\mathcal{H}_{\theta}\right)^{-1} m_{\theta}^{*}$ is well-defined. Note that the operator $\mathcal{H}_{\theta}$ is not necessarily a contraction, i.e., it may hold that $r\left(\mathcal{H}_{\theta}\right)>1$ for some $\theta \in \Theta$. Therefore, one cannot guarantee that the infinite sum $\sum_{j=0}^{\infty} \mathcal{H}_{\theta}^{j}$ exists for all $\theta \in \Theta$.

In practice one has to replace $m_{\theta}^{*}$ and $\mathcal{H}_{\theta}$ by estimators. Furthermore, one has also to estimate the parameters of the filters $A, B$. In the sequel we provide some details on this.

### 2.1.2 Further Details

Suppose we have a sample $\left\{\left(Y_{1}, X_{1}\right), \ldots,\left(Y_{T}, X_{T}\right)\right\}$. The general estimation strategy is

1. For each $\theta$ compute estimators of $\widehat{m}_{\theta}^{*}, \widehat{\mathcal{H}}_{\theta}$ of $m_{\theta}^{*}, \mathcal{H}_{\theta}$
2. Solve an empirical version of the equation (8) to obtain an estimator $\widehat{m}_{\theta}$ of $m_{\theta}$
3. Choose $\widehat{\theta}$ to minimize the profiled least squares criterion with respect to $\theta$. Let $\widehat{m}(x)=\widehat{m}_{\widehat{\theta}}(x)$.

Let $\tau=\tau(T)$ be some truncation parameter and define $Z_{t}^{\tau}(\beta)=\sum_{j=0}^{\tau} b_{j}(\beta) Y_{t-j}$. The choice of truncation depends on the dependence model $A(L), B(L)$. For geometrically declining parameters (as we shall assume) one can work with logarithmic truncation. There are many suitable estimators of the regression functions and density functions; we shall use local linear regression estimators for $m^{*}$ and a fairly standard kernel density estimator for $\mathcal{H}$ but other choices are possible.

For any sequence $\left\{Z_{t}^{\tau}(\beta)\right\}$ and any lag $j$ define the estimator $\widehat{g}_{j}(x ; \beta)=\widehat{c}_{0}$, where $\left(\widehat{c}_{0}, \widehat{c}_{1}\right)$ are the minimizers of the weighted sums of squares criterion

$$
\begin{equation*}
\sum_{t=j+1}^{T}\left\{Z_{t}^{\tau}(\beta)-c_{0}-c_{1}\left(X_{t-j}-x\right)\right\}^{2} K_{h}\left(X_{t-j}-x\right) \tag{12}
\end{equation*}
$$

with respect to $\left(c_{0}, c_{1}\right)$, where $K$ is a symmetric probability density function, $h$ is a positive bandwidth, and $K_{h}()=.K(. / h) / h$. Further define

$$
\begin{gathered}
\widehat{f}_{0, j}(y, x)=\frac{1}{T-|j|} \sum_{t=|j|+1}^{T} K_{h}\left(y, X_{t}\right) K_{h}\left(x, X_{t-j}\right) \quad ; \quad \widehat{f}_{0}(x)=\frac{1}{T} \sum_{t=1}^{T} K_{h}\left(x, X_{t}\right) . \\
\widehat{m}_{\theta}^{*}(x)=\sum_{j=0}^{\tau} a_{j}^{\dagger}(\alpha) \widehat{g}_{j}(x ; \beta) \quad ; \quad \widehat{\mathcal{H}}_{\theta}(x, y)=-\sum_{j= \pm 1}^{ \pm \tau} a_{j}^{*}(\alpha) \frac{\widehat{f}_{0, j}(y, x)}{\widehat{f}_{0}(y) \widehat{f}_{0}(x)} \\
\widehat{\mathcal{H}}_{\theta} m(x)=\int \widehat{\mathcal{H}}_{\theta}(x, y) m(y) \widehat{f}_{0}(y) d y .
\end{gathered}
$$

Here, 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 $X_{t}$. We shall assume that the covariate is supported on $[\underline{x}, \bar{x}]$ for some known $\underline{x}, \bar{x}$ and that the covariate density is bounded away from zero on this support. We need to make a boundary adjustment to the kernel $K$ in $\widehat{\mathcal{H}}_{\theta}$ by using the boundary kernels $K_{h}^{x}(y-x)$ to ensure that the bias is the same magnitude everywhere. One can also replace the standard kernel density estimators by other suitable density estimators like the Jones, Linton and Nielsen (1995) procedure, but it is not clear if such estimators would achieve better performance.

Then define $\widehat{m}_{\theta}$ as any solution to the equation

$$
\begin{equation*}
m=\widehat{m}_{\theta}^{*}+\widehat{\mathcal{H}}_{\theta} m \tag{13}
\end{equation*}
$$

in $L_{2}\left(\widehat{f}_{0}\right)$. We discuss the computation of this solution in the appendix. Let $\widehat{\theta}=\arg \min _{\theta \in \Theta} \widehat{Q}_{T}(\theta)$, where

$$
\widehat{Q}_{T}(\theta)=\frac{1}{T} \sum_{t=\tau+1}^{T}\left\{Z_{t}^{\tau}(\beta)-\sum_{j=0}^{\tau} a_{j}(\alpha) \widehat{m}_{\theta}\left(X_{t-j}\right)\right\}^{2}
$$

Finally, let $\widehat{m}(x)=\widehat{m}_{\widehat{\theta}}(x)$.

## 3 Asymptotic Properties

Let $\mathcal{F}_{a}^{b}$ be the $\sigma$-algebra of events generated by the random variables $\left\{Y_{t}, X_{t} ; a \leq j \leq b\right\}$. A stationary processes $\left\{Y_{t}, X_{t}\right\}$ is called strong mixing [Rosenblatt (1956)] if

$$
\begin{equation*}
\sup _{A \in \mathcal{F}_{-\infty}^{0}, B \in \mathcal{F}_{k}^{\infty}}|\operatorname{Pr}(A \cap B)-\operatorname{Pr}(A) \operatorname{Pr}(B)| \equiv s(k) \rightarrow 0 \quad \text { as } k \rightarrow \infty . \tag{14}
\end{equation*}
$$

We shall consider two cases. First, the 'weak form case' where we do not maintain that model (1) holds but only that $\left\{Y_{t}, X_{t}\right\}$ is a stationary strong mixing process. Second, we maintain that, in addition,

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model (1) holds with a martingale difference error sequence $\varepsilon_{t}$. To facilitate the asymptotic analysis, we make the following assumptions on the residuals and regressors, the kernel function $k(\cdot)$, and the bandwidth parameter $h$. Let $\eta_{t, j}(\beta)=Z_{t+j}(\beta)-E\left[Z_{t+j}(\beta) \mid X_{t}\right], \zeta_{t, j}(\theta)=m_{\theta}\left(X_{t+j}\right)-E\left[m_{\theta}\left(X_{t+j}\right) \mid X_{t}\right]$,

$$
\begin{equation*}
\eta_{\theta, t}^{1}=\sum_{j=0}^{\infty} a_{j}^{\dagger}(\alpha) \eta_{t, j}(\beta) \text { and } \eta_{\theta, t}^{2}=-\sum_{j= \pm 1}^{ \pm \infty} a_{j}^{*}(\alpha) \zeta_{t, j}(\theta) \tag{15}
\end{equation*}
$$

B1 The process $\left\{X_{t}, Y_{t}\right\}_{t=-\infty}^{\infty}$ is stationary and strong mixing with a mixing coefficient, $s(k)$ such that for some $C \geq 0$ and some $\bar{s}<1, s(k) \leq C \bar{s}^{k}$.

B2 $E\left(\left|Y_{t}\right|^{2 \rho}\right)<\infty$ for some $\rho>2$.
B3 The covariate process $\left\{X_{t}\right\}_{t=-\infty}^{\infty}$ has absolutely continuous density $f_{0}$ supported on $[\underline{x}, \bar{x}]$ for some $-\infty<\underline{x}<\bar{x}<\infty$ and the bivariate densities $f_{0, j}(\cdot)$ are supported on $[\underline{x}, \bar{x}]^{2}$. The function $m(\cdot)$ together with the densities $f_{0}(\cdot)$ and $f_{0, j}(\cdot)$ are continuous and twice continuously differentiable over $(\underline{x}, \bar{x})$ [and $(\underline{x}, \bar{x})^{2}$ ], and are uniformly bounded. $f_{0}(\cdot)$ is bounded away from zero on $[\underline{x}, \bar{x}]$, i.e., $\inf _{\underline{x} \leq w \leq \bar{x}} f_{0}(w)>0$.

B4 The parameter space $\Theta$ is a compact subset of $\mathbb{R}^{p}$, and the value $\theta_{0}$ is an interior point of $\Theta$. Also, A2 holds, and for any $\epsilon>0$

$$
\inf _{\left\|\theta-\theta_{0}\right\|>\epsilon} Q\left(\theta, m_{\theta}\right)>Q\left(\theta_{0}, m_{\theta_{0}}\right)
$$

B5 The density function $\mu$ of $\left(\eta_{t, j}^{1}(\beta), \eta_{t, j}^{2}(\beta)\right)$ is Lipschitz continuous on its domain. The joint densities $\mu_{0, j}, j=1,2, \ldots$, of $\left(\eta_{t, 0}^{1}(\beta), \eta_{t, 0}^{2}(\beta)\right),\left(\eta_{t, j}^{1}(\beta), \eta_{t, j}^{2}(\beta)\right)$ are uniformly bounded.

B6 The parameters $\alpha \in \mathcal{A}$ and $\beta \in \mathcal{B}$ compact subsets of $\mathbb{R}^{p_{a}}$ and $\mathbb{R}^{p_{b}}$ respectively. The coefficients satisfy $\sup _{\alpha \in A, k=0,1,2}\left\|\partial^{k} a_{j}(\alpha) / \partial \alpha^{k}\right\| \leq C \bar{a}^{j}$ for some $\bar{a}<1$ and some finite constant $C$, while $\inf _{\alpha \in \mathcal{A}} \sum_{j=0}^{\infty} a_{j}^{2}(\alpha)>0$. Likewise, $\sup _{\beta \in B, k=0,1,2}\left\|\partial^{k} b_{j}(\beta) / \partial \beta^{k}\right\| \leq C \bar{b}^{j}$ for some $\bar{b}<1$ and some finite constant $C$.

B7 The truncation sequence $\tau_{T}$ satisfies $\tau_{T}=C \log T$ for some constant $C>(-2 \log \bar{b})^{-1}$.
B8 The bandwidth sequence $h(T)$ satisfies $T^{1 / 5} h(T) \rightarrow \gamma$ as $T \rightarrow \infty$ with $\gamma$ bounded away from zero and infinity.

B9 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 $C, \sup _{x \in[x, \bar{x}]}\left|K^{x}(u)-K^{x}(v)\right| \leq C|u-v|$ for all $u, v \in[-1,1]$. Define $\mu_{j}(K)=\int u^{j} K(u) d u$ and $\|K\|_{2}^{2}=\int K^{2}(u) d u$.
$\mathrm{B} 10 \varepsilon_{t}$ satisfies $E\left[\varepsilon_{t} \mid\left\{X_{t-j}\right\}_{j=0}^{\infty},\left\{\varepsilon_{t-j}\right\}_{j=1}^{\infty}\right]=0$ a.s.
B11 (a) $\varepsilon_{t}$ is i.i.d. and independent of the process $\left\{X_{t}\right\}$; (b) $\varepsilon_{t}$ is also normally distributed.
These conditions are similar to Linton and Mammen (2005) but we also need conditions on the $b_{j}(\beta)$ coefficients and separate conditions on $X$ and $Y$.

Note that B1-B6 imply the uniform version of conditions A1-A2. Condition B1 rules out long memory but allows a wide range of processes used in practice. We will make use of the mixing property to apply the exponential inequality of Bosq (1998) and to establish a central limit theorem for $\widehat{m}_{\theta}$ in the weak form case. In this weak form case we can't apply martingale limit theory. We need to apply a central limit theorem to (local) averages of the processes $\eta_{\theta, t}^{1}$ and $\eta_{\theta, t}^{2}$ defined above. These processes need not be mixing but are near epoch dependent processes on the strong mixing bases $Y_{t}, X_{t}$ with exponentially declining weights under our conditions on $B$, $A$; we apply a CLT due to Lu (2001) for such processes using conditions B1 and B5,B6.

Condition B3 is quite standard in the nonparametric regression literature, Note that we only assume twice continuously differentiable $m$.

In B4 we explicitly assume the identification of the parametric part. We make this high level assumption for three reasons. First, we need identification in the weak case, and this seems like a natural assumption to make in view of our definition of the weak form process. Second, we allow the coefficients $a_{j}(\theta), b_{j}(\theta)$ to depend on $\theta$ in a complicated way. Third, the mapping $\theta \longmapsto m_{\theta}$ may be quite complicated to analyze. Hannan (1973) used high level conditions [c.f. his condition (4)] similar to ours.

The truncation rate assumed in B 7 is consistent with the exponential decaying mixing coefficients. It can be weakened at the expense of more detailed argumentation. In B8 we are anticipating a rate of convergence of $T^{-2 / 5}$ for $\widehat{m}_{\theta}$, which is consistent with second order smoothness on the function $m$. The assumptions B10 are expressed in terms of the unobserved $\left\{\varepsilon_{t-j}\right\}_{j=1}^{\infty}$ and are equivalent to assumptions on $\left\{y_{t-j}\right\}_{j=1}^{\infty}$ under an invertibility condition. Assumption B10 is needed for the consistency of the parameter estimates $\widehat{\theta}$. In the pure regression model $(2,3)$ one only needs a weaker assumption $E\left[\varepsilon_{t} \mid\left\{X_{t-j}\right\}_{j=0}^{\infty}\right]=0$ a.s. for consistent estimation of $m$ and $\theta$ as is known from the parametric case.

Define the functions $\beta_{\theta}=\left(I-\mathcal{H}_{\theta}\right)^{-1} \beta_{\theta}^{*}$ as solutions to the integral equations $\beta_{\theta}=\beta_{\theta}^{*}+\mathcal{H}_{\theta} \beta_{\theta}$, in which:

$$
\beta_{\theta}^{*}(x)=\frac{f_{0}^{\prime}(x)}{f_{0}(x)} \frac{\partial}{\partial x} \mathcal{H}_{\theta} m_{\theta}(x)+\mathcal{H}_{\theta} m_{\theta}^{\prime \prime}(x) .
$$

Then define

$$
\begin{aligned}
\omega_{\theta}(x) & =\frac{\|K\|_{2}^{2}}{f_{0}(x)} \operatorname{var}\left[\eta_{\theta, t}^{1}+\eta_{\theta, t}^{2}\right], \\
b_{\theta}(x) & =\frac{1}{2} \mu_{2}(K) \beta_{\theta}(x),
\end{aligned}
$$

where $\eta_{\theta, t}^{j}, j=1,2$ were defined above in (15). We prove the following theorem in the appendix.
Theorem 1. Suppose that B1-B9 hold. Then for each $\theta \in \Theta$ and $x \in(\underline{x}, \bar{x})$

$$
\begin{equation*}
\sqrt{T h}\left[\widehat{m}_{\theta}(x)-m_{\theta}(x)-h^{2} b_{\theta}(x)\right] \Longrightarrow N\left(0, \omega_{\theta}(x)\right), \tag{16}
\end{equation*}
$$

Both the bias and variance in this result are quite complicated even though a local linear smoother has been used in estimating $g_{j}$. This is a 'weak form' result, where the model (1) is not assumed.

We next maintain a 'semi-strong form' assumption B10, which requires the filters to be correctly specified. Under this assumption we can apply a CLT for martingale difference sequences. We obtain the properties of $\widehat{\theta}$ by an application of the asymptotic theory for semiparametric profiled estimators, see Severini and Wong (1992) and Newey (1994). This requires a uniform expansion for $\widehat{m}_{\theta}(x)$ and for the derivatives (with respect to $\theta$ ) of $\widehat{m}_{\theta}(x)$. Under B10, we get that

$$
\eta_{\theta_{0}, t}^{1}+\eta_{\theta_{0}, t}^{2}=\sum_{j=0}^{\infty} a_{j}^{\dagger} \varepsilon_{t+j} .
$$

Thus:

$$
\begin{equation*}
\omega_{\theta_{0}}(x)=\frac{\|K\|_{2}^{2} \sum_{j=0}^{\infty} a_{j}^{2}\left(\alpha_{0}\right) E\left[\varepsilon_{t}^{2} \mid X_{t-j}=x\right]}{f_{0}(x)\left[\sum_{j=0}^{\infty} a_{j}^{2}\left(\alpha_{0}\right)\right]^{2}} \tag{17}
\end{equation*}
$$

Let $\varepsilon_{t}(\theta)=Z_{t}(\beta)-\sum_{j=0}^{\infty} a_{j}(\alpha) m_{\theta}\left(X_{t-j}\right)$, and let

$$
\mathcal{J}=E\left[\frac{\partial^{2} \varepsilon_{t}}{\partial \theta \partial \theta^{\top}}\left(\theta_{0}\right)\right] \text { and } \mathcal{I}=E\left[\frac{\partial \varepsilon_{t}}{\partial \theta} \frac{\partial \varepsilon_{t}}{\partial \theta^{\top}} \varepsilon_{t}^{2}\left(\theta_{0}\right)\right] .
$$

Theorem 2. Suppose that Assumptions B1 to B10 hold. Then,

$$
\sqrt{T}\left(\widehat{\theta}-\theta_{0}\right) \Longrightarrow N\left(0, \mathcal{J}^{-1} \mathcal{I} \mathcal{J}^{-1}\right)
$$

Furthermore, for $x \in(\underline{x}, \bar{x})$

$$
\sqrt{T h}\left(\widehat{m}(x)-m(x)-h^{2} b_{\theta_{0}}(x)\right) \Longrightarrow N\left(0, \omega_{\theta_{0}}(x)\right) .
$$

## Remarks.

1. The quantities $\mathcal{J}, \mathcal{I}$, and $\omega_{\theta_{0}}(x)$ can be consistently estimated by a plug-in method thereby allowing inference, see Linton and Mammen (2005).
2. Under the 'strong form' assumption B11 $\widehat{\theta}$ is semiparametrically efficient, see Linton and Mammen (2005). There is generally an information loss from the necessity of estimating the function m.
3. For the nonparametric estimator the form of the asymptotic variance is quite natural as it decreases with the value of the covariate density and increases with the error variance. Under the 'strong form' special case $\mathrm{B} 11(\mathrm{a}), \omega_{\theta_{0}}(x)=\|K\|_{2}^{2} \sigma_{\varepsilon}^{2} / f_{0}(x) \sum_{j=0}^{\infty} a_{j}^{2}\left(\alpha_{0}\right)$. In the special case of nonparametric regression (2) and (3) we can compare the efficiency of our procedure with that of alternative estimators like the usual kernel regression estimator: this has asymptotic variance $\omega_{\text {Ker }}(x)=\|K\|_{2}^{2} \sigma_{\varepsilon}^{2} \sum_{j=0}^{\infty} c_{j}^{2} / f_{0}(x)$, where $C(L)=A(L)^{-1}$. We can also compare with the estimator of Xiao, Linton, Carroll, and Mammen (2003): this has asymptotic variance $\omega_{X L C M}(x)=\|K\|_{2}^{2} \sigma_{\varepsilon}^{2} / f_{0}(x)$. In this case,

$$
\begin{equation*}
\omega(x) \leq \omega_{X L C M}(x) \leq \omega_{K e r}(x) \tag{18}
\end{equation*}
$$

so that our procedure is always at least as efficient as these two competitors in terms of variance.
4. We now argue that our estimator is efficient in terms of variance for the special case that $\sigma_{j}^{2}(x)=E\left[\varepsilon_{t}^{2} \mid X_{t-j}=x\right]$ does not depend on $j$. Unfortunately, there is no Cramer-Rao lower bound to rely on in nonparametic estimation. We circumvent this problem by introducing a theoretical "oracle model" where more information is available and by showing that our estimator achieves the same asymptotic variance as a classical nonparametric estimator in the oracle model, see Linton (1997). We consider the following oracle model where one observes $X_{t}$ and $Z_{t j}$ for $1 \leq t \leq T, j \geq 0$ with

$$
Z_{t j}=a_{j} m\left(X_{t-j}\right)+\varepsilon_{t j}
$$

where $\left(\varepsilon_{t j}: 1 \leq t \leq T\right)$ are independent copies of $\left(\varepsilon_{t}: 1 \leq t \leq T\right)$. In this model each equation in the original model $Z_{t}=\sum_{j=0}^{\infty} a_{j} m\left(X_{t-j}\right)+\varepsilon_{t}$ is replaced by an infinite series of equations where for each new equation all summands beside one are known. Thus we argue that this model is more informative. In the oracle model for each value of $j$ a kernel estimator $\widehat{m}_{j}(x)$ of $m(x)$ is available. We now consider linear combinations $\widehat{m}_{w}(x)=\sum_{j=1}^{\infty} w_{j} \widehat{m}_{j}(x)$ with $\sum_{j=1}^{\infty} w_{j}=1$. This is the most efficient way to combine the information of $\widehat{m}_{j}(x)$ because these estimators are asymptotically independent. We now look for the weights $w_{j}$ that lead to minimal asymptotic variance. The variance is equal to $\sum_{j=1}^{\infty} w_{j}^{2} \sigma_{j}^{2}(x) a_{j}^{-2}\|K\|_{2}^{2} f_{0}(x)^{-1}$. This is minimized for the choice $w_{j}=a_{j}^{2} \sigma_{j}^{-2}(x) / \sum_{l=1}^{\infty} a_{l}^{2} \sigma_{l}^{-2}(x)$ resulting in the asymptotic variance

$$
\frac{\|K\|_{2}^{2} \sum_{j=1}^{\infty} a_{j}^{2} \sigma_{j}^{-2}(x)}{f_{0}(x)^{-1}\left[\sum_{j=1}^{\infty} a_{j}^{2} \sigma_{j}^{-2}(x)\right]^{2}} .
$$

This variance coincides with $\omega(x)$ if $\sigma_{j}^{2}(x)$ does not depend on $j$. If this is not the case one can proceed with a slight modification of our estimator that achieves the efficiency bound. Instead of
using (7) we use

$$
\sum_{j=0}^{\infty} b_{j} E\left[Z_{0} \mid X_{-j}=x\right]=\sum_{j=0}^{\infty} b_{j} a_{j} m(x)+\sum \sum_{j \neq k} a_{j} b_{k} E\left[m\left(X_{-j}\right) \mid X_{-k}=x\right]
$$

with $b_{j}=a_{j} \sigma_{j}^{-2}(x) / \sum_{l=1}^{\infty} a_{l}^{2} \sigma_{l}^{-2}(x)$. Proceeding as after equation (7) this results in another integral equation with different intercept $m^{*}$ and kernel $\mathcal{H}$. Because $b_{j}$ depends on $x$ both, $m^{*}$ and $\mathcal{H}$ depend also on $x$. That means for each value of $x$ one has to solve another integral equation. The empirical implementation of the method requires a consistent estimate of the conditional variances $\sigma_{j}^{2}(x)$.
5. The bias expression $b_{\theta_{0}}(x)$ does not apparently simplify, which makes comparison with other estimators according to mean squared error inconclusive. However, our estimator can be modified to get an asymptotic bias expression of the form that is asymptotically equivalent to

$$
\begin{equation*}
b(x)=\frac{1}{2} \mu_{2}(K) m^{\prime \prime}(x), \tag{19}
\end{equation*}
$$

which is as for a standard local linear estimator in regression. Then we get a straight mean squared error reduction over the local linear regression estimator. There exist two proposals for additive models to achieve a bias term of the form (19). The first approach is local linear smooth backfitting of Mammen, Linton, and Nielson (1999). In the smooth backfitting approach the backfitting algorithm updates estimates of the functions and its derivatives. This approach could be implemented in our setting by using an appropriate integral operator that acts on tuples of two functions (fitting $m$ and its derivative). A modification of local linear smooth backfitting has been proposed in Mammen and Park (2006). Their backfitting only uses one dimensional operators but achieves the same asymptotic behaviour as local linear smooth backfitting. In our setting their approach works as follows. One replaces $\widehat{\mathcal{H}}_{\theta}$ by

$$
\begin{gathered}
\widehat{\mathcal{H}}_{\theta}^{m o d}(y, x)=-\sum_{j= \pm 1}^{ \pm \tau_{T}} a_{j}^{*}(\alpha)\left[\frac{\widetilde{f}_{0 j}(y, x)}{\widetilde{f}_{0}(y) \widetilde{f}_{0}(x)}-\frac{\int \widetilde{f}_{0 j}(u, x) d u}{\widetilde{f}_{0}(x) \int \widetilde{f}_{0}(u) d u}\right] \\
\widetilde{f}_{0}(y)=\widehat{f}_{0}(y)-\frac{\widehat{f}_{0}^{*}(y)^{2}}{\widehat{f}_{0}^{* *}(y)} ; \quad \widetilde{f}_{0, j}(y, x)=\widehat{f}_{0, j}(y, x)-\frac{\widehat{f}_{0, j}^{*}(y, x) \widehat{f}_{0}^{*}(y)}{\widehat{f_{0}^{* *}}(y)}, \\
\widehat{f}_{0}^{*}(y)=\frac{1}{T-|j|} \sum_{t=|j|+1}^{T} K_{h}\left(X_{t}, y\right)\left(X_{t}-y\right) \quad ; \quad \widehat{f}_{0}^{* *}(y)=\frac{1}{T-|j|} \sum_{t=|j|+1}^{T} K_{h}\left(X_{t}, y\right)\left(X_{t}-y\right)^{2}, \\
\widehat{f_{0, j}}(y, x)=\frac{1}{T-|j|} \sum_{t=|j|+1}^{T} K_{h}\left(X_{t}, y\right) L_{h}\left(x, X_{t-j}\right) \quad ; \quad \widehat{f}_{0, j}^{*}(y, x)=\frac{1}{T-|j|} \sum_{t=|j|+1}^{T} K_{h}\left(X_{t}, y\right) L_{h}\left(x, X_{t-j}\right)\left(X_{t}-y\right) .
\end{gathered}
$$

Here, the kernel $L$ is defined as $L(u)=2 K_{1 / \sqrt{2}}(u)-K_{\sqrt{2}}(u)$. Furthermore, $L_{h}$ is defined as

$$
L_{h}(u, v)=\{a(v, h) u+b(v, h)\} L\left[h^{-1}(v-u)\right]
$$

with $a$ and $b$ chosen so that: $\int_{\underline{x}}^{\bar{x}} L_{h}(u, v) d u=1$ and $\int_{\underline{x}}^{\bar{x}}(u-v) L_{h}(u, v) d u=0$. Note that the integration runs over $u$ and not over $v$. Thus the kernel is not a boundary corrected kernel in the usual sense. A similar proposal has been made in Linton and Mammen (2005, p789).
6. The modified estimator has bias (19) and variance $\omega_{\theta_{0}}(x)$. The ranking in (18) then applies also to the comparison of mean squared errors. Our (modified) estimator is therefore better according to asymptotic mean squared error than these two estimators in this case.
7. The above efficiency discussion has just involved pairwise comparison of estimators. It may be possible to establish a further efficiency property following the work of Horowitz, Klemelä, and Mammen (2006). They showed in the context of additive nonparametric regression that the local linear smooth backfitting method of Mammen, Linton, and Nielsen (1999) is Best Linear Minimax, that is, there does not exist a linear procedure that has smaller minimax error.
8. The asymptotic distribution can be used to guide bandwidth selection. The IMSE optimal bandwidth is

$$
h=\left[\frac{\|K\|_{2}^{2}}{\mu_{2}^{2}(K)}\right]^{1 / 5}\left[\frac{\sigma_{\varepsilon}^{2}(\bar{x}-\underline{x})}{\sum_{j=0}^{\infty} a_{j}^{2}\left(\alpha_{0}\right) E\left[m^{\prime \prime}\left(X_{t}\right)^{2}\right]}\right]^{1 / 5} T^{-1 / 5}
$$

for the modified estimator under homoskedasticity, where $\sigma_{\varepsilon}^{2}$ is the variance of $\varepsilon_{t}$. In practice one must replace these quantities by estimates based on a parametric or nonparametric scheme.

## 4 A Nonstationary Case

In this section we investigate the case where $Y_{t}$ can be nonstationary but $X_{t}$ is stationary mixing as before. We wish to allow for the possibility of unit roots even if they might be quite rare in practical applications of this technology.

The most general case would be where both $A, B$ contained unit roots either simple or complex. For expositional reason we shall focus on the special case where $B(L)=A(L)=1-L$. Consider the model

$$
\begin{equation*}
(1-\rho L) Y_{t}=(1-\rho L) m\left(X_{t}\right)+\varepsilon_{t} \tag{20}
\end{equation*}
$$

where in fact $\rho_{0}=1$ and $\varepsilon_{t}$ obeys B11. In this case,

$$
\begin{equation*}
Y_{t}=m\left(X_{t}\right)+u_{t}, \tag{21}
\end{equation*}
$$

where $u_{t}=u_{t-1}+\varepsilon_{t}$ is a unit root process, Phillips (1987). We suppose that $u_{0}=0$.
Direct estimation of $Y_{t}$ on $X_{t}$ will produce inconsistent estimates of $m$. The Xiao, Linton, Carroll, and Mammen (2003) procedure is also inconsistent in this unit root case because it relies on the initial
standard nonparametric regression estimator that is inconsistent. On the other hand our estimation of the additive model

$$
Y_{t}-Y_{t-1}=m\left(X_{t}\right)-m\left(X_{t-1}\right)+\varepsilon_{t}
$$

with white noise errors will produce consistent estimates of $m$. In fact, the theory for $m_{\rho_{0}}$ is exactly as in Theorem 2. The task here is to determine that we can estimate the parameter $\rho$ in (20) consistently and thence estimate $m$ consistently. ${ }^{6}$

One issue is that for $\rho \neq 1$, the process $(1-\rho L) Y_{t}$ is non-stationary and so some of the definitions of the previous section do not make sense. Instead we define $m_{T \rho}$ to be the potentially time varying minimizer of

$$
Q_{T}(m)=\frac{1}{T} \sum_{t=1}^{T} E\left[\left\{Y_{t}-\rho Y_{t-1}-m\left(X_{t}\right)+\rho m\left(X_{t-1}\right)\right\}^{2}\right] .
$$

A necessary condition for $m$ to be the minimizer is that it satisfies the first order condition

$$
\begin{align*}
& \frac{1}{T} \sum_{t=1}^{T} E\left[Y_{t}-\rho Y_{t-1} \mid X_{t}=x\right]-\rho E\left[Y_{t}-\rho Y_{t-1} \mid X_{t-1}=x\right]  \tag{22}\\
& =\left(1+\rho^{2}\right) m_{T \rho}(x)-\rho\left(E\left[m_{T \rho}\left(X_{t}\right) \mid X_{t-1}=x\right]+E\left[m_{T \rho}\left(X_{t-1}\right) \mid X_{t}=x\right]\right)
\end{align*}
$$

Then note that from (21), $Y_{t}-\rho Y_{t-1}=m\left(X_{t}\right)-\rho m\left(X_{t-1}\right)+\varepsilon_{t}+(1-\rho) u_{t-1}$, and so $E\left[Y_{t}-\rho Y_{t-1} \mid X_{t}=x\right]$ and $E\left[Y_{t}-\rho Y_{t-1} \mid X_{t-1}=x\right]$ are time invariant. Furthermore, we have assumed that $X_{t}$ is stationary and $\varepsilon_{t}$ is i.i.d. Therefore, there exists a time invariant solution to equation (22). ${ }^{7}$ Furthermore, the solution is characterized by the integral equation (8) with in this special case:

$$
\begin{gathered}
m_{\rho}^{*}(x)=\frac{1}{1+\rho^{2}}\left(E\left[Y_{t}-\rho Y_{t-1} \mid X_{t}=x\right]-\rho E\left[Y_{t}-\rho Y_{t-1} \mid X_{t-1}=x\right]\right) \\
\mathcal{H}_{\rho}(x, y)=-\frac{\rho}{1+\rho^{2}}\left(\frac{f_{0,1}(y, x)}{f_{0}(y) f_{0}(x)}+\frac{f_{0,1}(x, y)}{f_{0}(y) f_{0}(x)}\right) .
\end{gathered}
$$

What is different here is the error in estimating $E\left[Y_{t}-\rho Y_{t-1} \mid X_{t-1}=x\right]$ for example can be large unless $\rho$ is close to one in which case the term $(1-\rho) u_{t-1}$ is small and the process $Y_{t}-\rho Y_{t-1}$ is almost stationary. The difference in behaviour of the resulting $\widehat{m}_{\rho}$ for $\rho=1$ and $\rho \neq 1$ is what drives the faster rate of convergence for $\widehat{\rho}$.

Define

$$
\widehat{Q}_{T}(\rho)=\frac{1}{T} \sum_{t=2}^{T}\left\{Y_{t}-\rho Y_{t-1}-\widehat{m}_{\rho}\left(X_{t}\right)+\rho \widehat{m}_{\rho}\left(X_{t-1}\right)\right\}^{2}
$$

[^4]and let $\hat{\rho}=\arg \min _{\rho} \widehat{Q}_{T}(\rho)$. We use a subset of the regularity conditions B that are relevant. Let $B$ denote the standard Brownian Motion on $[0,1]$.

Theorem 3. Suppose that assumption B1 holds for $X_{t}$, that B2 holds for $\varepsilon_{t}$, that B3, B7-B9 and B11 hold. Then

$$
T(\widehat{\rho}-1) \Longrightarrow \frac{\int_{0}^{1} B(s) d B(s)}{\int_{0}^{1} B^{2}(s) d s}
$$

Furthermore,

$$
\sqrt{T h}\left(\widehat{m}(x)-m(x)-h^{2} b(x)\right) \Longrightarrow N(0, \omega(x)),
$$

where

$$
b(x)=\left(I-\mathcal{H}_{1}\right)^{-1}\left[\frac{f_{0}^{\prime}}{f_{0}} \frac{\partial}{\partial x} \mathcal{H}_{1} m+\mathcal{H}_{1} m^{\prime \prime}\right](x)
$$

and

$$
\omega(x)=\|K\|_{2}^{2} \frac{E\left[\varepsilon_{t}^{2}\right]}{2 f_{0}(x)} .
$$

Note that the asymptotics for $\widehat{\rho}$ are the same as those of the infeasible least squares estimator $\bar{\rho}=\sum_{t=2}^{T} u_{t} u_{t-1} / \sum_{t=2}^{T} u_{t-1}^{2}$, so that estimation of $m$ has no effect on the limiting distribution. One can also obtain local to unity asymptotics which are the same as those of $\bar{\rho}$. The distribution theory can be used to perform a test of the null hypothesis of a unit root.

We remark that out theory can be generalized easily to allow for short run dynamics in addition to the unit root. Suppose that in (21), $(1-L) u_{t}=C(L) \varepsilon_{t}$, where $C(L)=\sum_{j=0}^{\infty} c_{j} L^{j}$ and $\sum_{j=0}^{\infty} j\left|c_{j}\right|<$ $\infty$. Then by the Beveridge-Nelson decomposition we have $u_{t}=C(1) \sum_{s=1}^{t} \varepsilon_{s}+C^{*}(L) \varepsilon_{t}$ under our assumptions, where $C^{*}(L)=\sum_{j=0}^{\infty} c_{j}^{*} L^{j}$ with $c_{j}^{*}=-\sum_{i=j+1}^{\infty} c_{i}$ being summable. Then the result in Theorem 3 follows (for the corresponding estimator) with the correction factor $C(1)$ in the variance.

## 5 Numerical Results

We investigate the performance of our procedure on simulated data in the context of a nonparametric regression with correlated errors. Our purpose is to compare the performance of our estimator to the natural competitor for that case, the local linear estimator. We focus on the relative performance of two optimally implemented methods to dispense with issues about bandwidth selection and the small sample performance of the benchmark estimator.

We suppose that

$$
Y_{t}=m\left(X_{t}\right)+u_{t}, u_{t}=\rho_{0} u_{t-1}+\varepsilon_{t}
$$

with $m(x)=\beta_{0} x^{2} / 2$, where $X_{t} \sim N(0,1)$, and $\varepsilon_{t} \sim N\left(0, \sigma_{\varepsilon}^{2}\right)$. We take $\beta_{0}=1$ and $\sigma_{\varepsilon}^{2}=1$. We examine the cases $T \in\{800,400,200\}$ and $\rho_{0} \in\{0,0.05,0.1, \ldots, 0.95,1.0\}$, and use $n s=1000$
replications. We compute our estimator $\widehat{m}$ using 200 grid points and use a grid search method to select $\rho \in\left[\rho_{0}-\delta, \rho_{0}+\delta\right]$ for $\delta=0.2 .{ }^{8}$ We also compute the standard local linear estimator $\widetilde{m}$, in both cases the Gaussian kernel was used.

We chose bandwidth to be optimal according to (asymptotic) weighted mean squared error

$$
P_{\infty}^{c}(\widehat{m})=\operatorname{plim}_{T \rightarrow \infty} T^{4 / 5} \int_{-c}^{c}[\widehat{m}(x)-m(x)]^{2} f_{0}(x) d x
$$

which gives $h_{\text {opt }}=c_{K} c_{M} T^{-1 / 5}$, where $c_{K}=\left(2 c\|K\|_{2}^{2} / \mu_{2}^{2}(K)\right)^{1 / 5}$ is to do with the kernel and $c_{M}=$ $\left(\sigma_{\varepsilon}^{2} /\left(1+\rho_{0}^{2}\right) \beta_{0}^{2}\left(F_{0}(c)-F_{0}(-c)\right)\right)^{1 / 5}$, where $F_{0}(x)$ is the c.d.f. of the covariate, is to do with the model. We have taken $c=2$, which corresponds to an interval containing almost $95 \%$ of the covariate distribution. For the standard local linear estimator the optimal bandwidth is $c_{K} c_{M}^{*} T^{-1 / 5}$ with $c_{M}^{*}=\left(\sigma_{\varepsilon}^{2} /\left(1-\rho_{0}^{2}\right) \beta_{0}^{2}\left(F_{0}(c)-F_{0}(-c)\right)^{1 / 5}\right.$ provided $\rho_{0} \neq 1$ (when $\rho_{0}=1$ we set $\rho_{0}$ in the formula arbitrarily to 0.95$)$.

In Figure 1 below we report the relative value of the performance measure $P_{T}(\widehat{m}) / P_{T}(\widetilde{m})$, where

$$
P_{T}(\widehat{m})=E \int_{-c}^{c}[\widehat{m}(x)-m(x)]^{2} f_{0}(x) d x
$$

and where $E$ is computed by the mean or median over Monte Carlo simulations and the integral is computed by the average over our grid of estimation points. ${ }^{9}$ Both estimators use their optimal bandwidths, and consequently their theoretical relative efficiency is $\left(\left(1-\rho_{0}^{2}\right) /\left(1+\rho_{0}^{2}\right)\right)^{4 / 5}$, which is independent of the other parameters. This is plotted below along with the simulation average value for the different sample sizes against $\rho$ values. The results indicate that $\widehat{m}$ is indeed more efficient than $\widetilde{m}$ and that the advantage takes off after $\rho_{0}=0.8$; until this value the advantage is less than $20 \%$ in MSE terms. For small values of $\rho_{0}$ the finite sample performance ratio is actually better than predicted, although this is partly because $\widetilde{m}$ performs worse than predicted by its asymptotic theory. Note that when $\rho_{0}=1$ the standard local linear estimator is inconsistent. The relative performance seems to get slightly worse with sample size. The absolute performance of both estimators improves with sample size but the MSE of $\widetilde{m}$ improves more rapidly in the relevant range of sample sizes than does the MSE of $\widehat{m}$.

[^5]

Figure 1. Shows the empirical performance ratio $P_{T}(\widehat{m}) / P_{T}(\widetilde{m})$ for different sample sizes along with the asymptotic value $P_{\infty}(\widehat{m}) / P_{\infty}(\widetilde{m})$ predicted from the asymptotic theory. $X_{t}$ iid $N(0,1)$.

We also looked at the case where $X_{t}$ is autocorrelated, specifically, $X_{t}=0.95 X_{t-1}+u_{t}$, where $u_{t}$ is normally distributed such that $X_{t}$ is marginally $N(0,1)$. Theoretically, this does not make any difference, and in practice if anything relative performance is improved for this case. The results are not shown for brevity.

We next examine the performance of $\widehat{\rho}$. When $\rho<1$ the MSE decreases pretty much as predicted and the distribution approximates a normal for the larger sample size. When $\rho_{0}=1$, our simulations show that the variance of $\widehat{\rho}$ decreases rapidly with sample size with standard deviation being 0.0161 , 0.00896 , and 0.00458 for $T=200,400$, and 800 respectively, which is consistent with superconsistency. Below we show the qq plots of the empirical quantiles against those of the Dicky-Fuller density in this unit root case. As the sample size increases the distribution approaches the asymptotic distribution.


Figure 2. Shows the q-q plots of $\widehat{\rho}$ against the Dicky-Fuller density for three different sample sizes.

$$
X_{t}=0.95 X_{t-1}+u_{t} \text { with } X_{t} \sim N(0,1)
$$

Overall these results are much better than obtained in Xiao et al. (2003) in terms of the small sample relative performance, and show in some cases substantial gains over the standard smoothing methods widely employed in practice. However, we acknowledge that in more complicated settings where the order of the polynomials $A, B$ is higher and perhaps has to be determined that the results will worsen.

## 6 Extensions

We conclude the paper with a discussion of two important extensions.

### 6.1 Nonstationary $X$

Suppose that

$$
X_{t}=X_{t-1}+\eta_{t}
$$

with $\eta_{t}$ also white noise and uncorrelated with $\varepsilon_{t}$. Thus $X_{t}$ is a unit root process. This makes a substantial difference to the asymptotics since the corresponding operator $\mathcal{H}_{\theta}(x, y)$ is random. Provided $X_{t}$ is null recurrent, we might expect consistency (Phillips and Park (1998)) but the rates of convergence are slower and the asymptotic distributions change. Simulation results support the consistency of $\widehat{m}$. In particular, the corresponding graphic to Figure 1 is almost identical.

An alternative type of nonstationarity for $X_{t}$ is deterministic trend. Suppose that

$$
\begin{equation*}
X_{t}=\mu(t / T)+\sigma(t / T) \eta_{t} \tag{23}
\end{equation*}
$$

where $\eta_{t}$ is a stationary mixing process, see Dahlhaus (1997). If $\sigma \equiv 0, X_{t}$ is purely deterministic. In this case, the asymptotics of kernel regression smoother are different and reflect the autocorrelation in $u_{t}$, see Hart (1991) and Fan and Yao (2003, Theorem 6.1). Also, there is a problem applying our method because of concurvity. Specifically, we have for any $j, m((t-j) / T)=m(t / T)+O(j / T)$ and so assumption A2 is violated. In this case we have $B(L) Y_{t} \simeq A(1) m(t / T)+\varepsilon_{t}$ and there appears to be no estimator that improves over the standard nonparametric regression. This is a bit like the well known result that OLS=GLS when the regressors are polynomial or trigonometric time trends. See Opsomer, Wang, and Yang (2001) for a review of nonparametric methods and results in this case. In the more general locally stationary case, our method may work due to the stochasticness of $\eta_{t}$.

### 6.2 Multivariate $X, Y$

When $X_{t}$ is multivariate the above method can be applied with obvious changes in the dimensionality of various quantities. However, it may be appealing in that case to consider the following model

$$
B(L) Y_{t}=\sum_{j=1}^{d} A_{j}(L) m_{j}\left(X_{j t}\right)+\varepsilon_{t}
$$

where the functions $m_{1}(),. \ldots, m_{d}($.$) are unknown and the filters A_{j}(L)=\sum_{k=0}^{\infty} a_{j k}, j=1, \ldots, d$. The estimation strategy involves a combination of Mammen, Linton, and Nielsen (1999) and the methods above. Instead one might want to make the function $m\left(X_{1 t}, \ldots, X_{d t}\right)$ obey some other dimensionality reducing restrictions.

## A Appendix

## A. 1 Computational Appendix

We discuss briefly how we solve the equation (13) in practice. Note that one can rewrite (9) as an integral equation on $[0,1]^{2}$ as $m_{\theta}^{\dagger}(s)=m_{\theta}^{* \dagger}(s)+\int_{0}^{1} \mathcal{H}_{\theta}^{\dagger}(s, t) m_{\theta}(t) d t$, where $\mathcal{H}_{\theta}^{\dagger}(s, t)=\mathcal{H}_{\theta}\left(F_{0}^{-1}(s), F_{0}^{-1}(t)\right)$
with $y=F_{0}^{-1}(s), x=F_{0}^{-1}(t)$ and $m_{\theta}^{\dagger}(t)=m_{\theta}\left(F_{0}^{-1}(t)\right)$ and $m_{\theta}^{* \dagger}(t)=m_{\theta}^{*}\left(F_{0}^{-1}(t)\right)$ and $F_{0}$ is the c.d.f. of $X_{t}$. For simplicity we drop the superfluous $\dagger$ superscript in the sequel. Let $\left\{t_{j, n}, j=1, \ldots, n\right\}$ be some equally spaced grid of points in $[0,1]$, and let $q_{j, n}=\widehat{F}_{0}^{-1}\left(t_{j, n}\right)$ be the empirical $t_{j, n}$ quantile of $X_{t}$. Now approximate (13) by

$$
\begin{equation*}
\widehat{m}_{\theta}\left(q_{i, n}\right)=\widehat{m}_{\theta}^{*}\left(q_{i, n}\right)+\sum_{j=1}^{n} \widehat{\mathcal{H}}_{\theta}\left(q_{i, n}, q_{j, n}\right) \widehat{m}_{\theta}\left(q_{j, n}\right), i=1, \ldots, n . \tag{24}
\end{equation*}
$$

The linear system (24) can be written in matrix notation

$$
\begin{equation*}
\left(I_{n}-\widehat{\mathbf{H}}_{\theta}\right) \widehat{\mathbf{m}}_{\theta}=\widehat{\mathbf{m}}_{\theta}^{*}, \tag{25}
\end{equation*}
$$

where $I_{n}$ is the $n \times n$ identity, $\widehat{\mathbf{m}}_{\theta}=\left(\widehat{m}_{\theta}\left(q_{1, n}\right), \ldots, \widehat{m}_{\theta}\left(q_{n, n}\right)\right)^{\top}$ and $\widehat{\mathbf{m}}_{\theta}^{*}=\left(\widehat{m}_{\theta}^{*}\left(q_{1, n}\right), \ldots, \widehat{m}_{\theta}^{*}\left(q_{n, n}\right)\right)^{\top}$, while

$$
\widehat{\mathbf{H}}_{\theta}=-\sum_{k= \pm 1}^{ \pm \tau} a_{k}^{+}(\alpha)\left[\frac{\widehat{f}_{0, k}\left(q_{i, n}, q_{j, n}\right)}{\widehat{f}_{0}\left(q_{i, n}\right) \widehat{f}_{0}\left(q_{j, n}\right)}\right]_{i, j=1}^{n}
$$

is an $n \times n$ matrix. We then find the solution values $\widehat{\mathbf{m}}_{\theta}=\left(\widehat{m}_{\theta}\left(q_{1, n}\right), \ldots, \widehat{m}_{\theta}\left(q_{n, n}\right)\right)^{\top}$ to this system (25) by direct inversion when $n$ is less than say 2000. The error that is made in such discrete approximations is well studied, see for example Porter and Stirling (1990). Under our smoothness conditions, we expect the error to decline like $O(1 / n)$ as $n \rightarrow \infty$.

We next describe a method proposed to us by a referee, which follows from arguments of Carrasco, Florens, and Renault (2006). This method actually provides an explicit solution to the integral equation (13). Write the operator action on function $m$ as

$$
\left(\widehat{\mathcal{H}}_{\theta} m\right)(x)=-\sum_{j= \pm 1}^{ \pm \tau} a_{j}^{*}(\alpha) \frac{1}{\widehat{f}_{0}(x)} \frac{1}{T-|j|} \sum_{t=|j|+1}^{T} K_{h}\left(x, X_{t-j}\right) \xi\left(X_{t}\right),
$$

where $\xi\left(X_{t}\right)=\int K_{h}\left(y, X_{t}\right) m(y) d y$. Premultiplying both sides of the integral equation (13) by $K_{h}\left(x, X_{t}\right)$ and integrating over $x$, we obtain a system of $T$ linear equations

$$
\begin{equation*}
\widehat{\xi}_{\theta}=\widehat{U}_{\theta} \widehat{\xi}_{\theta}+\widehat{w}_{\theta}, \tag{26}
\end{equation*}
$$

where $\widehat{\xi}_{\theta}$ is the $T \times 1$ vector with typical element $\widehat{\xi}_{\theta}\left(X_{t}\right)=\int K_{h}\left(y, X_{t}\right) \widehat{m}_{\theta}(y) d y, \widehat{U}_{\theta}=\sum_{j= \pm 1}^{ \pm \tau} \widehat{U}_{\theta j}$, where $\widehat{U}_{\theta j}$ is the $T \times T$ matrix with $(t, t-j)$ element equal to

$$
-\sum_{j= \pm 1}^{ \pm \tau} a_{j}^{*}(\alpha) \frac{1}{\widehat{f}_{0}(x)} \frac{1}{T-|j|} \int K_{h}\left(x, X_{t}\right) K_{h}\left(x, X_{t-j}\right) d x
$$

if $t \geq|j|+1$ and zero otherwise, and $\widehat{w}_{\theta}$ is the $T \times 1$ vector with typical element $\int K_{h}\left(x, X_{t}\right) \widehat{m}_{\theta}^{*}(x) d x$. Provided $I-\widehat{U}_{\theta}$ is invertible, $\widehat{\xi}_{\theta}=\left(I-\widehat{U}_{\theta}\right)^{-1} \widehat{w}_{\theta}$. One then obtains $\widehat{m}_{\theta}(x)$ from the original equation
$m(x)=\widehat{m}_{\theta}^{*}(x)+\left(\widehat{\mathcal{H}}_{\theta} m\right)(x)$. The integrals $\int K_{h}\left(x, X_{t}\right) K_{h}\left(x, X_{t-j}\right) d x$ can be computed explicitly in some cases. This method is a nice complement to the method we have used in practice. The advantage of the method is that it provides an exact solution. The disadvantage is that it can be quite computationally demanding since one has to invert $T \times T$ matrices, which obviously limits the sample sizes one can work with. In addition, the above method only yields a solution at a single point $x$, whereas the method we have used automatically yields a solution at all grid points.

## A. 2 Proof of Theorems

## A.2.1 Stationary Case

Proof of Theorem 1. The proof strategy uses the general result in Linton and Mammen (2005) for the treatment of empirical integral equations. First, for general $\theta$ we apply Linton and Mammen (2005, Proposition 1). Thus we write

$$
\begin{align*}
\widehat{m}_{\theta}^{*}(x)-m_{\theta}^{*}(x) & =\widehat{m}_{\theta}^{*, B}(x)+\widehat{m}_{\theta}^{*, C}(x)+\widehat{m}_{\theta}^{*, D}(x)  \tag{27}\\
\left(\widehat{\mathcal{H}}_{\theta}-\mathcal{H}_{\theta}\right) m_{\theta}(x) & =\widehat{m}_{\theta}^{*, E}(x)+\widehat{m}_{\theta}^{*, F}(x)+\widehat{m}_{\theta}^{*, G}(x), \tag{28}
\end{align*}
$$

where $\widehat{m}_{\theta}^{*, B}(x)$ and $\widehat{m}_{\theta}^{*, E}(x)$ are deterministic and $O\left(T^{-2 / 5}\right)$,

$$
\begin{aligned}
& \widehat{m}_{\theta}^{*, B}(x)=\frac{h^{2}}{2} \mu_{2}(K)\left[m_{\theta}^{\prime \prime}(x)-\frac{\partial^{2}}{\partial x^{2}}\left(\mathcal{H}_{\theta} m_{\theta}\right)(x)\right] \\
& \widehat{m}_{\theta}^{*, E}(x)=\frac{h^{2}}{2} \mu_{2}(K)\left[\mathcal{H}_{\theta} m_{\theta}^{\prime \prime}(x)+\frac{\partial^{2}}{\partial x^{2}}\left(\mathcal{H}_{\theta} m_{\theta}\right)(x)+2 \frac{f_{0}^{\prime}(x)}{f_{0}(x)} \frac{\partial}{\partial x}\left(\mathcal{H}_{\theta} m_{\theta}\right)(x)\right]
\end{aligned}
$$

while:

$$
\begin{aligned}
\widehat{m}_{\theta}^{*, C}(x) & =\frac{1}{T f_{0}(x)} \sum_{t} K_{h}\left(x, X_{t}\right) \eta_{\theta, t}^{1} \\
\widehat{m}_{\theta}^{*, F}(x) & =\frac{1}{T f_{0}(x)} \sum_{t} K_{h}\left(x, X_{t}\right) \eta_{\theta, t}^{2}
\end{aligned}
$$

and the remainder terms $\widehat{m}_{\theta}^{*, D}(x)$ and $\widehat{m}_{\theta}^{*, G}(x)$ satisfy

$$
\begin{array}{ll}
\sup _{\theta \in \Theta} \sup _{x \in \mathcal{X}_{\text {int }}}\left|\widehat{m}_{\theta}^{*, j}(x)\right|=o_{p}\left(T^{-2 / 5}\right), & j=D, G \\
\sup _{\theta \in \Theta} \sup _{x \in[\underline{x}, \bar{x}]}\left|\widehat{m}_{\theta}^{*, j}(x)\right|=O_{p}\left(T^{-2 / 5}\right), & j=D, G
\end{array}
$$

where $\mathcal{X}_{\text {int }}=[\underline{x}+h, \bar{x}-h]$.
From this one obtains an expansion

$$
\begin{equation*}
\widehat{m}_{\theta}(x)-m_{\theta}(x)-\left[m_{\theta}^{B}(x)+m_{\theta}^{E}(x)+\widehat{m}_{\theta}^{*, C}(x)+\widehat{m}_{\theta}^{*, F}(x)\right]=o_{p}\left(T^{-2 / 5}\right), \tag{29}
\end{equation*}
$$

uniformly for $\theta \in \Theta$ and for $x \in \mathcal{X}_{\text {int }}$ and $=O_{p}\left(T^{-2 / 5}\right)$, elsewhere. Here $m_{\theta}^{B}=\left(I-\mathcal{H}_{\theta}\right)^{-1} \widehat{m}_{\theta}^{*, B}$ and $m_{\theta}^{E}=\left(I-\mathcal{H}_{\theta}\right)^{-1} \widehat{m}_{\theta}^{*, E}$. From this expansion we obtain the main result. Specifically, $\widehat{m}_{\theta}^{*, C}(x)+\widehat{m}_{\theta}^{*, F}(x)$ is asymptotically normal with zero mean and the stated variance after applying a CLT for near epoch dependent functions of mixing processes. The asymptotic bias comes from $m_{\theta}^{B}(x)+m_{\theta}^{E}(x)$. Note that because of the boundary modification to the kernel we have $E \widehat{f}_{0}(x)=f_{0}(x)+O\left(h^{2}\right)$ and $E \widehat{f}_{0, j}(x, y)=f_{0, j}(x, y)+O\left(h^{2}\right)$ for all $x, y$.

Our proof below make use of the following results. For $\delta_{T}=T^{-3 / 10+\xi}$ with $\xi>0$ small enough,

$$
\begin{align*}
\max _{1 \leq j \leq \tau_{T}} \sup _{x, y \in[\underline{x}, \bar{x}]}\left|\widehat{f}_{0, j}(x, y)-f_{0, j}(x, y)\right| & =o_{p}\left(\delta_{T}\right)  \tag{30}\\
\sup _{x \in[\underline{x}, \bar{x}]}\left|\widehat{f}_{0}(x)-f_{0}(x)\right| & =o_{p}\left(\delta_{T}\right) \tag{31}
\end{align*}
$$

This follows by the exponential inequality of Bosq (1998, Theorem 1.3), see Linton and Mammen (2005, p817).

Proof of (27). Write

$$
Z_{t}(\beta)-Z_{t}^{\tau}(\beta)=\sum_{j=\tau+1}^{\infty} b_{j}(\beta) Y_{t-j} .
$$

We have $E\left[Z_{t}(\beta)-Z_{t}^{\tau}(\beta)\right]=E\left[Y_{t}\right] \sum_{j=\tau+1}^{\infty} b_{j}(\beta)=O\left(\bar{b}^{\tau}\right)$ and

$$
\begin{aligned}
\operatorname{var}\left[Z_{t}(\beta)-Z_{t}^{\tau}(\beta)\right] & =\sum_{j=\tau+1}^{\infty} \sum_{j^{\prime}=\tau+1}^{\infty} b_{j}(\beta) b_{j^{\prime}}(\beta) \operatorname{cov}\left(Y_{t-j}, Y_{t-j^{\prime}}\right) \\
& \leq \sum_{j=\tau+1}^{\infty} \sum_{j^{\prime}=\tau+1}^{\infty}\left|b_{j}(\beta)\right|\left|b_{j^{\prime}}(\beta)\right|\left|\gamma_{Y}\left(j-j^{\prime}\right)\right| \\
& \leq \sup _{u}\left\langle\gamma_{Y}(u)\right|\left(\sum_{j=\tau+1}^{\infty}\left|b_{j}(\beta)\right|\right)^{2}=O\left(\bar{b}^{2 \tau}\right)=o\left(T^{-1}\right)
\end{aligned}
$$

for each $\beta$. Similar bounds can be obtained for the covariance $\operatorname{cov}\left[Z_{t}(\beta)-Z_{t}^{\tau}(\beta), Z_{s}(\beta)-Z_{s}^{\tau}(\beta)\right]$. Let $\widetilde{g}_{j}(x ; \beta)$ denote (12) with $Z_{t}(\beta)$ replacing $Z_{t}^{\tau}(\beta)$. Then

$$
\begin{equation*}
\max _{1 \leq j \leq \tau} \sup _{\beta \in \mathcal{B}} \sup _{x \in[\underline{x}, \bar{x}]}\left|\widehat{g}_{j}(x ; \beta)-\widetilde{g}_{j}(x ; \beta)\right|=o_{p}\left(T^{-1 / 2}\right) . \tag{32}
\end{equation*}
$$

This follows using the above moment bounds and because of the assumed uniform decay rates on $b_{j}(\beta)$ and its derivatives and the moment condition on $Y$. See Xiao et al. (2003) for a similar argument.

Then for each $j \geq 0$,

$$
\widetilde{g}_{j}(x ; \beta)-g_{j}(x ; \beta)=\frac{1}{T f_{0}(x)} \sum_{t=1}^{T-j} K_{h}\left(x, X_{t}\right) \eta_{t, j}(\beta)+\frac{h^{2}}{2} \mu_{2}(K) \mathbf{b}_{j}(x ; \beta)+R_{T j}(x ; \beta),
$$

where $\mathbf{b}_{j}(x ; \beta)$ is the bias function and $R_{T j}(x ; \beta)$ is the remainder term, which is $o_{p}\left(T^{-2 / 5}\right)$ uniformly over $j \leq \tau_{T}, x \in[\underline{x}, \bar{x}]$ and $\beta \in \mathcal{B}$. By interchanging the order of summation we obtain for $x \in \mathcal{X}_{\text {int }}$

$$
\begin{aligned}
\sum_{j=0}^{\tau} a_{j}^{\dagger}(\alpha) \sum_{t=1}^{T-j} K_{h}\left(x, X_{t}\right) \eta_{t, j}(\beta) & =\sum_{t=1}^{T-\tau} \sum_{j=0}^{T-t} K_{h}\left(x, X_{t}\right) a_{j}^{\dagger}(\alpha) \eta_{t, j}(\beta) \\
& =\sum_{t=1}^{T-\tau} K_{h}\left(x, X_{t}\right) \sum_{j=0}^{\infty} a_{j}^{\dagger}(\alpha) \eta_{t, j}(\beta) \\
& -\sum_{t=1}^{T-\tau} K_{h}\left(x, X_{t}\right) \sum_{j=T-t+1}^{\infty} a_{j}^{\dagger}(\alpha) \eta_{t, j}(\beta) \\
= & \sum_{t=1}^{T-\tau} K_{h}\left(x, X_{t}\right) \sum_{j=0}^{\infty} a_{j}^{\dagger}(\alpha) \eta_{t, j}(\beta) \\
& -\sum_{j=\tau+1}^{T} a_{j}^{\dagger}(\alpha) \sum_{t=T-j+1}^{T-\tau} K_{h}\left(x, X_{t}\right) \eta_{t, j}(\beta) \\
& -\sum_{j=T+1}^{\infty} a_{j}^{\dagger}(\alpha) \sum_{t=1}^{T-\tau} K_{h}\left(x, X_{t}\right) \eta_{t, j}(\beta)
\end{aligned}
$$

where the terms apart from the first are of smaller order. Specifically,

$$
\begin{align*}
& \max _{1 \leq j \leq \tau} \sup _{\theta \in \Theta} \sup _{x \in[x, \bar{x}]}\left|\frac{1}{T f_{0}(x)} \sum_{j=\tau+1}^{T} a_{j}^{\dagger}(\alpha) \sum_{t=T-j+1}^{T-\tau} K_{h}\left(x, X_{t}\right) \eta_{t, j}(\beta)\right|=o_{p}\left(T^{-2 / 5}\right)  \tag{33}\\
& \max _{1 \leq j \leq \tau} \sup _{\theta \in \Theta} \sup _{x \in[x, \bar{x}]}\left|\frac{1}{T f_{0}(x)} \sum_{j=T+1}^{\infty} a_{j}^{\dagger}(\alpha) \sum_{t=1}^{T-\tau} K_{h}\left(x, X_{t}\right) \eta_{t, j}(\beta)\right|=o_{p}\left(T^{-2 / 5}\right) . \tag{34}
\end{align*}
$$

These follow by standard arguments. Therefore,

$$
\begin{aligned}
\sum_{j=0}^{\tau} a_{j}^{\dagger}(\alpha)\left[\widehat{g}_{j}(x ; \beta)-g_{j}(x ; \beta)\right] & =\frac{1}{T f_{0}(x)} \sum_{t=1}^{T-\tau} K_{h}\left(x, X_{t}\right) \sum_{j=0}^{\infty} a_{j}^{\dagger}(\alpha) \eta_{t, j}(\beta) \\
& +\frac{h^{2}}{2} \mu_{2}(K) \sum_{j=0}^{\tau} a_{j}^{\dagger}(\alpha) \mathbf{b}_{j}(x ; \beta)+o_{p}\left(T^{-2 / 5}\right)
\end{aligned}
$$

uniformly over $x \in[\underline{x}, \bar{x}]$. Claim (27) now follows from
$\widehat{m}_{\theta}^{*, B}(x)=\frac{h^{2}}{2} \mu_{2}(K) \sum_{j=0}^{\tau} a_{j}^{\dagger}(\alpha) \mathbf{b}_{j}(x ; \beta)+o_{p}\left(T^{-2 / 5}\right)$ uniformly for $x \in \mathcal{X}_{\text {int }}$ and $\theta \in \Theta$,
$\widehat{m}_{\theta}^{*, B}(x)=\frac{h^{2}}{2} \mu_{2}(K) \sum_{j=0}^{\tau} a_{j}^{\dagger}(\alpha) \mathbf{b}_{j}(x ; \beta)+O_{p}\left(T^{-2 / 5}\right)$ uniformly for $x \in x \in[\underline{x}, \bar{x}] \backslash \mathcal{X}_{\text {int }}$ and $\theta \in \Theta$.

Proof of (28). We have

$$
\begin{aligned}
& \int \widehat{\mathcal{H}}_{\theta}(x, y) m_{\theta}(y) \widehat{f}_{0}(y) d y-\int \mathcal{H}_{\theta}(x, y) m_{\theta}(y) f_{0}(y) d y \\
& =-\sum_{j= \pm 1}^{ \pm \tau} a_{j}^{*}(\alpha) \int\left[\frac{\widehat{f}_{0, j}(x, y)}{\widehat{f_{0}}(x)}-\frac{f_{0, j}(x, y)}{f_{0}(x)}\right] m_{\theta}(y) d y+\sum_{j= \pm \tau \pm 1}^{ \pm \infty} a_{j}^{*}(\alpha) \int \frac{f_{0, j}(x, y)}{f_{0}(x)} m_{\theta}(y) d y \\
& =-\sum_{j= \pm 1}^{ \pm \tau} a_{j}^{*}(\alpha) \int\left[\frac{\widehat{f}_{0, j}(x, y)}{\widehat{f}_{0}(x)}-\frac{f_{0, j}(x, y)}{f_{0}(x)}\right] m_{\theta}(y) d y+o\left(T^{-2 / 5}\right)
\end{aligned}
$$

uniformly over $x, \theta$ due to the uniform decay rates on $a_{j}(\alpha)$. Specifically,

$$
\sup _{\theta \in \Theta} \sup _{x \in[x, \bar{x}]}\left|\sum_{j= \pm \tau \pm 1}^{ \pm \infty} a_{j}^{*}(\alpha) \int \frac{f_{0, j}(x, y)}{f_{0}(x)} m_{\theta}(y) d y\right| \leq C \bar{a}^{\tau} \times \bar{m}=o\left(T^{-2 / 5}\right),
$$

where $\sup _{\theta \in \Theta} \sup _{y \in[x, \bar{x}]}\left|m_{\theta}(y)\right|=\bar{m}<\infty$.
Denote by

$$
\int \frac{f_{0, j}(x, y)}{f_{0}(x)} m_{\theta}(y) d y=E\left[m\left(X_{t-j}\right) \mid X_{t}=x\right] \equiv r_{j}(x) .
$$

Then write

$$
\begin{equation*}
\frac{\int \widehat{f}_{0, j}(x, y) m_{\theta}(y) d y}{\widehat{f}_{0}(x)}=\frac{\frac{1}{T} \sum_{t} K_{h}\left(x, X_{t}\right) m_{t-j}^{*}}{\frac{1}{T} \sum_{t} K_{h}\left(x, X_{t}\right)} \tag{35}
\end{equation*}
$$

where

$$
m_{t}^{*}=\int K_{h}^{y}\left(y-X_{t}\right) m_{\theta}(y) d y
$$

Then note that for $X_{t} \in \mathcal{X}_{\text {int }}$

$$
\begin{aligned}
\int K_{h}^{y}\left(y-X_{t}\right) m_{\theta}(y) d y-m_{\theta}\left(X_{t}\right) & =\int K_{h}^{y}\left(y-X_{t}\right)\left[m_{\theta}(y)-m_{\theta}\left(X_{t}\right)\right] d y \\
& =m_{\theta}^{\prime}\left(X_{t}\right) \int K_{h}^{y}\left(y-X_{t}\right)\left(y-X_{t}\right) d y \\
& +\frac{1}{2} \int K_{h}^{y}\left(y-X_{t}\right)\left(y-X_{t}\right)^{2} m_{\theta}^{\prime \prime}\left(X_{t}^{*}(y)\right) d y \\
& =\frac{h^{2}}{2} \mu_{2}(K) m_{\theta}^{\prime \prime}\left(X_{t}\right)+o\left(h^{2}\right)
\end{aligned}
$$

by a second order Taylor expansion, a change of variables and property B9 of the kernels. The error is uniformly $o\left(h^{2}\right)$ over $t$ with $X_{t} \in \mathcal{X}_{\text {int }}$ and $\theta$. Note that (35) is just like a local constant smoother of $m_{t-j}^{*}$ on $X_{t}$ and can be analyzed in the same way.

Using $\widehat{a} / \widehat{b}-c=(\widehat{a}-\widehat{b} c) / \widehat{b}$, we have

$$
\begin{align*}
& \frac{\int \widehat{f}_{0, j}(x, y) m_{\theta}(y) d y}{\widehat{f}_{0}(x)}-\int \frac{f_{0, j}(x, y)}{f_{0}(x)} m_{\theta}(y) d y \\
& =\frac{\frac{1}{T} \sum_{t} K_{h}\left(x, X_{t}\right)\left[m_{t-j}^{*}-r_{j}(x)\right]}{\frac{1}{T} \sum_{t} K_{h}\left(x, X_{t}\right)} \\
& =\frac{\frac{1}{T} \sum_{t} K_{h}\left(x, X_{t}\right)\left[m_{\theta}\left(X_{t-j}\right)-r_{j}(x)\right]}{\frac{1}{T} \sum_{t} K_{h}\left(x, X_{t}\right)}+\frac{\frac{1}{T} \sum_{t} K_{h}\left(x, X_{t}\right)\left[m_{t-j}^{*}-m_{\theta}\left(X_{t-j}\right)\right]}{\frac{1}{T} \sum_{t} K_{h}\left(x, X_{t}\right)}  \tag{36}\\
& \simeq \frac{\frac{1}{T} \sum_{t} K_{h}\left(x, X_{t}\right)\left[m_{\theta}\left(X_{t-j}\right)-r_{j}\left(X_{t}\right)\right]}{\frac{1}{T} \sum_{t} K_{h}\left(x, X_{t}\right)}+\frac{\frac{1}{T h} \sum_{t} K_{h}\left(x, X_{t}\right)\left[r_{j}\left(X_{t}\right)-r_{j}(x)\right]}{\frac{1}{T h} \sum_{t} K_{h}\left(x, X_{t}\right)} \\
& +\frac{h^{2}}{2} \mu_{2}(K) E\left[m_{\theta}^{\prime \prime}\left(X_{t-j}\right) \mid X_{t}=x\right] \\
& \simeq \frac{1}{T h} \frac{1}{f_{0}(x)} \sum_{t} K_{h}\left(x-X_{t}\right) \zeta_{t, j}+\frac{h^{2}}{2} \mu_{2}(K)\left[r_{j}^{\prime \prime}(x)+\frac{2 r_{j}^{\prime}(x) f_{0}^{\prime}(x)}{f_{0}(x)}+E\left[m_{\theta}^{\prime \prime}\left(X_{t-j}\right) \mid X_{t}=x\right]\right] \tag{37}
\end{align*}
$$

by standard arguments for Nadaraya-Watson smoothers. The approximation is of order $o\left(T^{-2 / 5}\right)$, uniformly over $j \leq \tau_{T}$, over $x$ in $\mathcal{X}_{\text {int }}$, and over $\theta \in \Theta$. Summing this up, gives (28) for $x \in \mathcal{X}_{\text {int }}$. The proof for the boundary follows by standard arguments.

Proof of Theorem 2. The consistency of $\hat{\theta}$ follows along the lines of Linton and Mammen (2005) using the expansions obtained above uniform over $\theta$. Note that the solution value $m_{\theta}$ is twice continuously differentiable in $\theta$ under our assumptions and

$$
\begin{align*}
\frac{\partial m_{\theta}}{\partial \theta} & =\left(\frac{\partial m_{\theta}^{*}}{\partial \theta}+\frac{\partial \mathcal{H}_{\theta}}{\partial \theta} m_{\theta}\right)+\mathcal{H}_{\theta} \frac{\partial m_{\theta}}{\partial \theta}  \tag{38}\\
\frac{\partial^{2} m_{\theta}}{\partial \theta \partial \theta^{\top}} & =\left(\frac{\partial^{2} m_{\theta}^{*}}{\partial \theta \partial \theta^{\top}}+\frac{\partial^{2} \mathcal{H}_{\theta}}{\partial \theta \partial \theta^{\top}} m_{\theta}+\frac{\partial \mathcal{H}_{\theta}}{\partial \theta} \frac{\partial m_{\theta}}{\partial \theta^{\top}}\right)+\mathcal{H}_{\theta} \frac{\partial^{2} m_{\theta}}{\partial \theta \partial \theta^{\top}} . \tag{39}
\end{align*}
$$

These define $\partial m_{\theta} / \partial \theta$ and $\partial^{2} m_{\theta} / \partial \theta \partial \theta^{\top}$ as solutions to integral equations with different intercepts but the same operator $\mathcal{H}_{\theta}$ as (10), so the solution to these equations exists and is unique by the arguments given above.

Let $Q(\theta)=Q\left(\theta, m_{\theta}\right)$ with $Q\left(\theta, m_{\theta}\right)$ defined in (11). We first show that

$$
\begin{equation*}
\sup _{\theta \in \Theta}\left|\widehat{Q}_{T}(\theta)-Q(\theta)\right| \xrightarrow{P} 0 \tag{40}
\end{equation*}
$$

which follows from $\sup _{\theta \in \Theta} \sup _{x \in[\underline{x}, \bar{x}]}\left|\widehat{m}_{\theta}(x)-m_{\theta}(x)\right| \xrightarrow{P} 0$ given the moment and mixing conditions etc. This follows from the expansions in Theorem 1 and standard uniform convergence arguments for
kernel smoothers. Specifically, $\sup _{\theta \in \Theta} \sup _{x \in[\underline{x}, \bar{x}]}\left|\widehat{m}_{\theta}^{*, j}(x)\right|=o_{p}(1), \quad j=B, C$. The uniformity over $\theta$ comes from analysis of $\partial m_{\theta}^{*}(x) / \partial \theta$ and $\partial \widehat{m}_{\theta}^{*, j}(x) / \partial \theta$. Then apply assumption B 4 to yield consistency of $\widehat{\theta}$.

Define the score function and Hessian

$$
\begin{aligned}
\frac{\partial \widehat{Q}_{T}(\theta)}{\partial \theta} & =\frac{1}{T} \sum_{t=2}^{T} \widehat{\varepsilon}_{t}^{\tau}(\theta) \frac{\partial \widehat{\varepsilon}_{t}^{\tau}(\theta)}{\partial \theta} \\
\frac{\partial^{2} \widehat{Q}_{T}(\theta)}{\partial \theta \partial \theta^{\top}} & =\frac{1}{T} \sum_{t=2}^{T} \frac{\partial \widehat{\varepsilon}_{t}^{\tau}(\theta)}{\partial \theta} \frac{\partial \widehat{\varepsilon}_{t}^{\tau}(\theta)}{\partial \theta^{\top}}+\widehat{\varepsilon}_{t}^{\tau}(\theta) \frac{\partial^{2} \widehat{\varepsilon}_{t}^{\tau}(\theta)}{\partial \theta \partial \theta^{\top}}
\end{aligned}
$$

where $\widehat{\varepsilon}_{t}^{\tau}(\theta)=Z_{t}^{\tau}(\beta)-\sum_{j=0}^{\tau} a_{j}(\alpha) \widehat{m}_{\theta}\left(X_{t-j}\right)$. One then establishes a CLT for the score function at $\theta=\theta_{0}$ and a local uniform law of large numbers for the Hessian, which establish the CLT for $\widehat{\theta}$.

We can now effectively take $\theta=\theta_{0}$ in Theorem 1 . The asymptotic statement on the distribution of $\widehat{m}_{\theta_{0}}(x)-m(x)$ directly follows from Theorem 1 . Note that $\eta_{\theta, t}^{1}+\eta_{\theta, t}^{2}=\sum_{j=0}^{\infty} a_{j}^{*} \varepsilon_{t+j}$.

## A.2.2 Nonstationary Case

Proof of Theorem 3. Let

$$
\begin{gathered}
\varepsilon_{t}(\rho)=Y_{t}-\rho Y_{t-1}-m_{\rho}\left(X_{t}\right)+\rho m_{\rho}\left(X_{t-1}\right)=Y_{t}-\rho Y_{t-1}-m\left(X_{t}\right)+\rho m\left(X_{t-1}\right)=\varepsilon_{t}+(1-\rho) u_{t-1} \\
\widehat{\varepsilon}_{t}(\rho)=Y_{t}-\rho Y_{t-1}-\widehat{m}_{\rho}\left(X_{t}\right)+\rho \widehat{m}_{\rho}\left(X_{t-1}\right)
\end{gathered}
$$

We have

$$
\begin{aligned}
Q_{T}(\rho) & =\frac{1}{T} \sum_{t=2}^{T} \varepsilon_{t}^{2}(\rho)=\frac{1}{T} \sum_{t=2}^{T} \varepsilon_{t}^{2}+T(1-\rho)^{2} \frac{1}{T^{2}} \sum_{t=2}^{T} u_{t-1}^{2}+2(1-\rho) \frac{1}{T} \sum_{t=2}^{T} \varepsilon_{t} u_{t-1} \\
& \simeq \sigma_{\varepsilon}^{2}+T(1-\rho)^{2} \sigma_{\varepsilon}^{2} \int B^{2}(s) d s+2(1-\rho) \sigma_{\varepsilon}^{2} \int B(s) d B(s)
\end{aligned}
$$

The least squares estimator that minimizes $Q_{T}(\rho)$, denoted $\bar{\rho}$, has closed form expression $\bar{\rho}=$ $\sum_{t=2}^{T} u_{t} u_{t-1} / \sum_{t=2}^{T} u_{t-1}^{2}$. It is consistent at rate $T$ and furthermore

$$
\begin{equation*}
T(\bar{\rho}-1) \Longrightarrow \frac{\int B(s) d B(s)}{\int B^{2}(s) d s} \tag{41}
\end{equation*}
$$

We next consider the difference between $\widehat{Q}_{T}(\rho)$ and $Q_{T}(\rho)$. We have

$$
\begin{gather*}
\widehat{Q}_{T}(\rho)=Q_{T}(\rho)+\frac{1}{T} \sum_{t=2}^{T}\left\{\widehat{\varepsilon}_{t}(\rho)-\varepsilon_{t}(\rho)\right\}^{2}+2 \frac{1}{T} \sum_{t=2}^{T}\left\{\widehat{\varepsilon}_{t}(\rho)-\varepsilon_{t}(\rho)\right\} \varepsilon_{t}(\rho),  \tag{42}\\
\widehat{\varepsilon}_{t}(\rho)-\varepsilon_{t}(\rho)=-\left(\widehat{m}_{\rho}\left(X_{t}\right)-m_{\rho}\left(X_{t}\right)\right)+\rho\left(\widehat{m}_{\rho}\left(X_{t-1}\right)-m_{\rho}\left(X_{t-1}\right)\right) .
\end{gather*}
$$

Proof of Consistency. We prove:

$$
\begin{equation*}
\widehat{Q}_{T}(1) \rightarrow^{p} q \tag{43}
\end{equation*}
$$

for some $q>0$ (hence $\widehat{Q}_{T}(1) / T \rightarrow^{p} 0$ ), and

$$
\begin{equation*}
\lim _{T \rightarrow \infty} \inf _{|\rho-1|>\delta} \frac{1}{T} \widehat{Q}_{T}(\rho)>0 \tag{44}
\end{equation*}
$$

Combine (43) and (44) yields $\widehat{\rho} \xrightarrow{P} 1$.
Proof of (43). The properties of $\widehat{Q}_{T}(1)$ can be derived using the expansion of Theorem 1, and specifically the uniform over $x$ consistency of $\widehat{m}_{1}(x)$. We have $\widehat{Q}_{T}(1) \xrightarrow{P} E\left(\varepsilon_{t}^{2}\right)>0$.

Proof of (44). We first derive the properties of $\widehat{m}_{\rho}-m_{\rho}$ for $\rho \neq 1$. As in the stationary case we can approximate $\widehat{m}_{\rho}-m_{\rho}$ in terms of $\widehat{m}_{\rho}^{*}-m_{\rho}^{*}$ and $\left(\widehat{\mathcal{H}}_{\rho}-\mathcal{H}_{\rho}\right) m_{\rho}$. The expansion for $\left(\widehat{\mathcal{H}}_{\rho}-\mathcal{H}_{\rho}\right) m_{\rho}$ is as above. The main difference concerns the fact that the expansion for $\widehat{m}_{\rho}^{*}-m_{\rho}^{*}$ contains a term that is large when $\rho \neq 1$ and indeed $\widehat{m}_{\rho}^{*}$ does not consistently estimate $m_{\rho}^{*}$ unless $\rho=1$. Therefore, $\widehat{m}_{\rho}-m_{\rho}$ is dominated by the large term in $\widehat{m}_{\rho}^{*}-m_{\rho}^{*}$. Specifically, (28) holds but (27) needs to be modified.

The intercept function $m_{\rho}^{*}$ is

$$
m_{\rho}^{*}(x)=\frac{1}{1+\rho^{2}}\left(E\left[Y_{t}-\rho Y_{t-1} \mid X_{t}=x\right]-\rho E\left[Y_{t}-\rho Y_{t-1} \mid X_{t-1}=x\right]\right)=\frac{1}{1+\rho^{2}}\left[g_{0 \rho}(x)-\rho g_{1 \rho}(x)\right]
$$

a linear combination of $g_{0 \rho}(x)=E\left[Y_{t}-\rho Y_{t-1} \mid X_{t}=x\right]$ and $g_{1 \rho}(x)=E\left[Y_{t}-\rho Y_{t-1} \mid X_{t-1}=x\right]$. Therefore, we must establish the properties of $\widehat{g}_{j \rho}(x)-g_{j \rho}(x), j=0,1$, where $\widehat{g}_{j \rho}(x)$ are the estimates of $g_{j \rho}(x)$ when $\rho \neq 1$. We have
$Y_{t}-\rho Y_{t-1}-E\left[Y_{t}-\rho Y_{t-1} \mid X_{t}=x\right]=m\left(X_{t}\right)-m(x)-\rho\left(m\left(X_{t-1}\right)-E\left[m\left(X_{t-1}\right) \mid X_{t}=x\right]\right)+\varepsilon_{t}+(1-\rho) u_{t-1}$.
$Y_{t}-\rho Y_{t-1}-E\left[Y_{t}-\rho Y_{t-1} \mid X_{t-1}=x\right]=m\left(X_{t}\right)-E\left[m\left(X_{t}\right) \mid X_{t-1}=x\right]-\rho\left(m\left(X_{t-1}\right)-m(x)\right)+\varepsilon_{t}+(1-\rho) u_{t-1}$.
The terms $m\left(X_{t}\right)-m(x)$ and $m\left(X_{t-1}\right)-m(x)$ on the right hand side contribute to biases; the stationary error terms $-\rho\left(m\left(X_{t-1}\right)-E\left[m\left(X_{t-1}\right) \mid X_{t}=x\right]\right)+\varepsilon_{t}$ and $m\left(X_{t}\right)-E\left[m\left(X_{t}\right) \mid X_{t-1}=x\right]+\varepsilon_{t}$ may contribute to the variance but are standard, it is the term $(1-\rho) u_{t-1}$ containing the unit root that is different. We have

$$
\begin{aligned}
\widehat{g}_{j \rho}(x)-g_{j \rho}(x) & =\frac{1}{T f_{0}(x)} \sum_{t=j+1}^{T} K_{h}\left(x, X_{t-j}\right) \varepsilon_{t}+(1-\rho) \frac{1}{T f_{0}(x)} \sum_{t=j+1}^{T} K_{h}\left(x, X_{t-j}\right) u_{t-1} \\
& +\frac{h^{2}}{2} \mu_{2}(K) \mathbf{b}_{j}(x ; \rho)+R_{T}(x ; \rho) \equiv \delta_{T 1}(x)+\delta_{T 2}(x)+\delta_{T 3}(x)+R_{T}(x ; \rho),
\end{aligned}
$$

where $\sup _{x \in[x, \bar{x}]} \delta_{T 1}(x)=O_{p}\left(\sqrt{\log T} T^{-2 / 5}\right)$ and $\sup _{x \in[x, \bar{x}]} \delta_{T 3}(x)=O_{p}\left(T^{-2 / 5}\right)$ under our bandwidth conditions, while the remainder term is of smaller order than $\delta_{T 2}(x)$. This approximation is valid because the $X$ process is stationary so the terms except $\delta_{T 2}(x)$ are standard.

We consider the term $\delta_{T 2}(x)$ and write $\delta_{T 2}(x)=\sqrt{T}(1-\rho) \xi_{T}(x)+\sqrt{T}(1-\rho) \eta_{T}(x)$ with

$$
\begin{aligned}
& \xi_{T}(x)=\frac{1}{T} \sum_{t=1}^{T} E\left[\frac{1}{f_{0}(x)} K_{h}\left(x, X_{t-j}\right)\right] \frac{u_{t-1}}{\sqrt{T}} \\
& \eta_{T}(x)=\frac{1}{T} \sum_{t=1}^{T}\left(\frac{1}{f_{0}(x)} K_{h}\left(x, X_{t-j}\right)-E\left[\frac{1}{f_{0}(x)} K_{h}\left(x, X_{t-j}\right)\right]\right) \frac{u_{t-1}}{\sqrt{T}}
\end{aligned}
$$

Clearly, because $E\left[\frac{1}{f_{0}(x)} K_{h}\left(x, X_{t-j}\right)\right]=1+O\left(h^{2}\right)$ uniformly in $x$,

$$
\xi_{T}(x)=\frac{1}{T} \sum_{t=1}^{T} \frac{u_{t-1}}{\sqrt{T}}+o_{p}(1)=O_{p}(1)
$$

uniformly in $x$.
We argue that $\sup _{x \in[\underline{x}, \bar{x}]}\left|\eta_{T}(x)\right|=o_{p}(1)$. Note that $E\left[\eta_{T}(x)\right]=0$ by assumption B11. Define

$$
\begin{equation*}
\epsilon_{T t}=\frac{1}{f_{0}(x)} K_{h}\left(x, X_{t-j}\right)-E\left[\frac{1}{f_{0}(x)} K_{h}\left(x, X_{t-j}\right)\right] . \tag{45}
\end{equation*}
$$

This has (approximately as $T \rightarrow \infty$ ) covariance function

$$
\begin{aligned}
\operatorname{cov}\left(\epsilon_{T t}, \epsilon_{T t-r}\right) & =E\left[\frac{1}{f_{0}^{2}(x)} K_{h}\left(x, X_{t}\right) K_{h}\left(x, X_{t-r}\right)\right]-E^{2}\left[\frac{1}{f_{0}(x)} K_{h}\left(x, X_{t}\right)\right] \\
& \simeq \frac{f_{0, t-r}(x, x)}{f_{0}^{2}(x)}-1 \equiv \gamma_{\epsilon}(t-r),
\end{aligned}
$$

by the standard change of variable and dominated convergence argument.
Furthermore,

$$
\begin{aligned}
\operatorname{var}\left[\eta_{T}(x)\right] & =\frac{1}{T^{3}} \sum_{t=j+1}^{T} E\left[\epsilon_{T t}^{2}\right] E\left[u_{t-1}^{2}\right]+\frac{1}{T^{3}} \sum_{t \neq s} E\left[\epsilon_{T t} \epsilon_{T s}\right] E\left[u_{t} u_{s}\right] \\
& \simeq \frac{\sigma_{\varepsilon}^{2}}{T^{3}} \sum_{t \neq s} \sum_{t=s} \min \{s, t\} \gamma_{\epsilon}(t-s) \simeq \frac{2 \sigma_{\varepsilon}^{2}}{T^{3}} \sum_{s=1}^{T-1} s \sum_{t=s+1}^{T} \gamma_{\epsilon}(t-s) \\
& \simeq \frac{2 \sigma_{\varepsilon}^{2}}{T^{2}} \sum_{s=1}^{T-1} s(T-s) \sum_{k=1}^{\infty} \gamma_{\epsilon}(k)=\frac{2 \sigma_{\varepsilon}^{2}}{3 T} \sum_{k=1}^{\infty} \gamma_{\epsilon}(k),
\end{aligned}
$$

so that $\operatorname{var}\left[\eta_{T}(x)\right]=O\left(T^{-1}\right)$ and $\eta_{T}(x)=O_{p}\left(T^{-1 / 2}\right)$ for each $x \in[\underline{x}, \bar{x}]$. The pointwise result can be extended to uniformity over $x \in[\underline{x}, \bar{x}]$ by standard arguments, so $\sup _{x \in[\underline{x}, \bar{x}]}\left|\eta_{T}(x)\right|=o_{p}(1)$ as required. Therefore

$$
\begin{equation*}
\widehat{g}_{j \rho}(x)-g_{j \rho}(x)=\sqrt{T}(1-\rho) \frac{1}{T} \sum_{t=1}^{T} \frac{u_{t-1}}{\sqrt{T}}+o_{p}(\sqrt{T}) \tag{46}
\end{equation*}
$$

Note that the right hand side is the same regardless of location $x$ and $j$ and the error is uniform over these quantities. By the usual arguments (Phillips (1987)), $T^{-3 / 2} \sum_{t=1}^{T} u_{t-1} \Longrightarrow \sigma_{\varepsilon} \int_{0}^{1} B(s) d s$. Therefore, $\left(\widehat{g}_{j \rho}(x)-g_{j \rho}(x)\right) / \sqrt{T} \Longrightarrow(1-\rho) \sigma_{\varepsilon} \int_{0}^{1} B(s) d s$ for all $x$ and $j=0,1$. In fact this convergence is uniform over $x$.

It holds that:

$$
\begin{gather*}
\frac{1}{T^{2}} \sum_{t=2}^{T}\left\{\widehat{\varepsilon}_{t}(\rho)-\varepsilon_{t}(\rho)\right\}^{2}=\frac{(1-\rho)^{6}}{\left(1+\rho^{2}\right)^{2}} \sigma_{\varepsilon}^{2}\left(\int_{0}^{1} B(s) d s\right)^{2}+o_{p}(1)  \tag{47}\\
(1-\rho) \frac{1}{T^{2}} \sum_{t=2}^{T}\left\{\widehat{\varepsilon}_{t}(\rho)-\varepsilon_{t}(\rho)\right\} u_{t-1}=\frac{-(1-\rho)^{4}}{1+\rho^{2}} \sigma_{\varepsilon}^{2}\left(\int_{0}^{1} B(s) d s\right)^{2}+o_{p}(1)  \tag{48}\\
\frac{1}{T} \sum_{t=2}^{T}\left\{\widehat{\varepsilon}_{t}(\rho)-\varepsilon_{t}(\rho)\right\} \varepsilon_{t}=\frac{-(1-\rho)^{3}}{1+\rho^{2}} \frac{1}{\sqrt{T}} \sum_{t=2}^{T} \frac{u_{t-1}}{\sqrt{T}} \varepsilon_{t}+o_{p}(1)=O_{p}(1) . \tag{49}
\end{gather*}
$$

We just show the argument for (47). We have

$$
\begin{aligned}
& \frac{1}{T^{2}} \sum_{t=2}^{T}\left\{\widehat{\varepsilon}_{t}(\rho)-\varepsilon_{t}(\rho)\right\}^{2} \\
& =\frac{1}{T^{2}} \sum_{t=2}^{T}\left\{\left(\widehat{m}_{\rho}\left(X_{t}\right)-m\left(X_{t}\right)\right)-\rho\left(\widehat{m}_{\rho}\left(X_{t-1}\right)-m\left(X_{t-1}\right)\right)\right\}^{2} \\
& =\frac{1}{T^{2}} \sum_{t=2}^{T}\left\{\left(\widehat{m}_{\rho}^{*}-m_{\rho}^{*}\right)\left(X_{t}\right)-\rho\left(\widehat{m}_{\rho}^{*}-m_{\rho}^{*}\right)\left(X_{t-1}\right)\right\}^{2}+o_{p}(1) \\
& =\frac{1}{\left(1+\rho^{2}\right)^{2}} \frac{1}{T^{2}} \sum_{t=2}^{T}\left\{\left[\widehat{g}_{0 \rho}-g_{0 \rho}\right]\left(X_{t}\right)+\rho^{2}\left[\widehat{g}_{1 \rho}-g_{1 \rho}\right]\left(X_{t-1}\right)-\rho\left[\widehat{g}_{0 \rho}-g_{0 \rho}\right]\left(X_{t-1}\right)-\rho\left[\widehat{g}_{1 \rho}-g_{1 \rho}\right]\left(X_{t}\right)\right\}^{2} \\
& +o_{p}(1) \\
& =\frac{(1-\rho)^{6}}{\left(1+\rho^{2}\right)^{2}}\left\{\frac{1}{T} \sum_{t=1}^{T} \frac{u_{t-1}}{\sqrt{T}}\right\}^{2}+o_{p}(1)
\end{aligned}
$$

by (46). From this (47) follows. The arguments for (48) and (49) are similar.

Then, by (47)-(49) we have

$$
\begin{aligned}
\widehat{Q}_{T}(\rho) & \simeq \sigma_{\varepsilon}^{2}+T(1-\rho)^{2} \sigma_{\varepsilon}^{2} \int_{0}^{1} B^{2}(s) d s+2(1-\rho) \sigma_{\varepsilon}^{2} \int_{0}^{1} B(s) d B(s) \\
& +\frac{(1-\rho)^{6} T}{\left(1+\rho^{2}\right)^{2}} \sigma_{\varepsilon}^{2}\left(\int_{0}^{1} B(s) d s\right)^{2}-\frac{2(1-\rho)^{4} T}{1+\rho^{2}} \sigma_{\varepsilon}^{2}\left(\int_{0}^{1} B(s) d s\right)^{2}-\frac{-2(1-\rho)^{3}}{1+\rho^{2}} O_{p}(1) .
\end{aligned}
$$

Therefore

$$
\begin{aligned}
\frac{1}{T} \widehat{Q}_{T}(\rho) & \simeq(1-\rho)^{2} \sigma_{\varepsilon}^{2} \int_{0}^{1} B^{2}(s) d s+\left[\frac{(1-\rho)^{6}}{\left(1+\rho^{2}\right)^{2}}-\frac{2(1-\rho)^{4}}{1+\rho^{2}}\right] \sigma_{\varepsilon}^{2}\left(\int_{0}^{1} B(s) d s\right)^{2} \\
& =(1-\rho)^{2} \sigma_{\varepsilon}^{2} \int_{0}^{1} B^{2}(s) d s-\frac{(1-\rho)^{4}(\rho+1)^{2}}{\left(1+\rho^{2}\right)^{2}} \sigma_{\varepsilon}^{2}\left(\int_{0}^{1} B(s) d s\right)^{2}
\end{aligned}
$$

By the Cauchy-Schwarz inequality $\int_{0}^{1} B^{2}(s) d s \geq\left(\int_{0}^{1} B(s) d s\right)^{2}$. Therefore, with probability one:

$$
\begin{equation*}
\frac{1}{T} \widehat{Q}_{T}(\rho) \geq 4(1-\rho)^{2} \frac{\rho^{2}}{\left(1+\rho^{2}\right)^{2}} \sigma_{\varepsilon}^{2}\left(\int_{0}^{1} B(s) d s\right)^{2}>0 \tag{50}
\end{equation*}
$$

for all $\rho \neq 1$. This establishes (44).
Proof of Asymptotic Distribution. Reparameterizing $\rho \mapsto r=1-\rho / T$ we get

$$
\widehat{Q}_{T}(r) \simeq \sigma_{\varepsilon}^{2}+\frac{r^{2}}{T} \sigma_{\varepsilon}^{2} \int_{0}^{1} B^{2}(s) d s+2 \frac{r}{T} \sigma_{\varepsilon}^{2} \int_{0}^{1} B(s) d B(s)+o\left(T^{-1}\right)
$$

so that the terms from the nonparametric estimation drop out. Therefore, the asymptotic distribution is the Dickey-Fuller distribution, i.e.,

$$
T(\hat{\rho}-1) \Longrightarrow \frac{\int_{0}^{1} B(s) d B(s)}{\int_{0}^{1} B^{2}(s) d s}
$$

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[^1]:    ${ }^{1}$ Sims (1971) and Geweke (1978) consider a continuous time distributed lag model where $Y(t)=\int_{-\infty}^{\infty} a(s) X(t-$ $s) d s+\varepsilon(t)$ and the data are observed at discrete time intervals in which case the (high frequency) discrete time approximation to this is like (1) with $B(L)=1$ and $A(L)=\sum_{j=-\infty}^{\infty} a_{j} L^{j}$ for some $a_{j}$ related to the function $a($.) under some conditions.
    ${ }^{2}$ The usual properties of linear dynamic regression models can be extended to the nonlinear case. Thus for example we can define the average instantaneous impact $E\left[\partial Y_{t} / \partial X_{t}\right]$ as equal to the average derivative of the function $m$, = $E\left[m^{\prime}\left(X_{t}\right)\right]$, a quantity that has been investigated elsewhere. The total dynamic average impact $\sum_{j=0}^{\infty} E\left[\partial Y_{t+j} / \partial X_{t}\right]=$ $E\left[m^{\prime}\left(X_{t}\right)\right] \sum_{j=0}^{\infty}(B(L) / A(L))_{j}$ is proportional to the instantaneous impact.

[^2]:    ${ }^{3}$ This equation can also be derived at by directly taking conditional expectations of $Z_{t}$ given each $X_{t-k}$, multiplying by $a_{k}$, and then summing over $k$.

[^3]:    ${ }^{4}$ These are real numbers for which there exists functions $e_{j}($.$) such that \mathcal{H} e_{j}=\lambda_{j} e_{j}$.
    ${ }^{5}$ One example where this condition is not satisfied (and where the subsequent conclusions also do not hold) is when $X_{t}=t / T$.

[^4]:    ${ }^{6}$ Differencing can be expected to eliminate unit roots so long as enough differencing is undertaken. However, differencing produces additive models for which the optimal estimation strategy is a similar type of method to ours.
    ${ }^{7}$ Note also that $m_{\rho}=m$ for all $\rho$.

[^5]:    ${ }^{8}$ The value $\delta=0.2$ was large enough in our case such that nearly all values of $\widehat{\rho}$ were interior points of the parameter set.
    ${ }^{9}$ We also examined an integrate mean absolute error performance measure, but the results were similar and are not reported here.

