# Dynamic time series binary choice* 

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#### Abstract

This paper considers dynamic time series binary choice models. It proves near epoch dependence and strong mixing for the dynamic binary choice model with correlated errors. Using this result, it shows in a time series setting the validity of the dynamic probit likelihood procedure when lags of the dependent binary variable are used as regressors, and it establishes the asymptotic validity of Horowitz' smoothed maximum score estimation of dynamic binary choice models with lags of the dependent variable as regressors. For the semiparametric model, the latent error is explicitly allowed to be correlated. It turns out that no long-run variance estimator is needed for the validity of the smoothed maximum score procedure in the dynamic time series framework.


## 1 Introduction

For a dynamic linear time series model

$$
\begin{equation*}
y_{n}=\sum_{j=1}^{p} \rho_{j} y_{n-j}+\gamma^{\prime} x_{n}+u_{n}, \tag{1}
\end{equation*}
$$

$n=1, \ldots, N$, it is well-known that a sufficient condition for consistency as $N \rightarrow \infty$ of the least squares estimator is that $E\left(u_{n} \mid y_{n-1}, \ldots, y_{n-p}, x_{n}\right)=0$, and that even if $u_{n}$ is weakly dependent, consistency can be proven as long as this condition holds, without the assumption

[^0]of normality on $u_{n}$. In this paper, we analyze maximum likelihood estimation of the dynamic probit model of order $p$, and maximum score estimation of dynamic binary choice models of order $p$, and we explicitly allow the error to be correlated. We define the dynamic binary choice model of order $p$ as
\[

$$
\begin{equation*}
y_{n}=I\left(\sum_{j=1}^{p} \rho_{j} y_{n-j}+\gamma^{\prime} x_{n}+u_{n}>0\right) \tag{2}
\end{equation*}
$$

\]

where $x_{n}$ is predetermined and $u_{n}$ can be correlated and heteroskedastic. We first show near epoch dependence and strong mixing for this model. We then impose identifying assumptions to ensure identification of the probit model and the binary choice model. For conditional maximum likelihood estimation of the dynamic probit model, the key condition that is needed will turn out to be

$$
\begin{equation*}
E\left(y_{n} \mid x_{n}, y_{n-1}, y_{n-2}, \ldots\right)=\Phi\left(\sum_{j=1}^{p} \rho_{j} y_{n-j}+\gamma^{\prime} x_{n}\right) \tag{3}
\end{equation*}
$$

while in the smoothed maximum score setting, we will need the condition

$$
\begin{equation*}
\operatorname{Median}\left(u_{n} \mid y_{n-1}, \ldots, y_{n-p}, x_{n}\right)=0 \tag{4}
\end{equation*}
$$

Therefore, this paper analyzes the dynamic time series binary choice model at a level of generality that is comparable to the level of generality at which linear dynamic time series models can be analyzed.
Manski (1975) uses the sign function to develop the first semiparametric estimator for the binary choice model. Cosslett (1983) and Ichimura (1993) derive alternative estimators for the binary choice model. Imbens (1992) and Matzkin (1992) also develop estimators for the semiparametric binary choice model. Finally, in his seminal paper, Horowitz (1992) smooths the sign function of Manski $(1975,1985)$ and derives an estimator that is asymptotically normally distributed. However, all these estimators assume that one has a random sample. Thus, none of these estimators allows for lagged dependent explanatory variables. Park and Phillips (2000) assume that one of the regressors in a binary choice model is integrated, and they assume that all regressors are exogenous, thereby excluding predetermined variables and lagged $y_{n}$ as possible regressors. Other recent papers that consider multinomial choice models in the presence of an integrated regressor are Hu and Phillips (2004) and Moon (2004).

In this paper we consider the binary choice model in a time series setting and we allow for lagged dependent variables and predetermined regressors as explanatory variables. For the semiparametric case, we only impose a median assumption. Thus, we allow the variance (and
other moments of the error distribution) to depend on lagged error terms, lagged dependent variables as well as regressors. Moreover, the median assumption allows for heterogeneity that is caused by random coefficients, e.g. a data generating process whose parameters are random and symmetrically distributed around $\left(\rho^{\prime}, \gamma^{\prime}\right)^{\prime}$.
Ruud (1981) and Poirier and Ruud (1988) have considered the probit model with correlated errors. Robinson (1982) considered the tobit model with correlated errors. An example of an empirical paper that focuses on the dynamic probit methodology is Eichengreen, Watson, and Grossman (1985). However, no formal stationarity properties for dynamic probit models are derived in these papers, nor anywhere else in the literature as far as the authors are aware. Potential applications include finance models concerning the likelihood of a financial transaction in a given time period as well as models concerning labor market participation decisions in which the relative importance of wealth versus welfare effects are studied.
The setup of this paper is as follows. In Section 2, the weak dependence properties of $y_{n}$ are analyzed. Section 3 of this paper will analyze the dynamic probit procedure when lagged values of $y_{n}$ have been included among the regressors and normality of $u_{n}$ is assumed. In Section 4, we consider consistency of the smoothed maximum score estimator of the dynamic time series binary choice model. The smoothed maximum score estimator was first suggested in Horowitz (1992). Section 5 establishes asymptotic normality of the smoothed maximum score estimator ${ }^{1}$.

## 2 Properties of the dynamic time series binary choice model

A key aspect of the analysis below is to show that $y_{n}$ satisfies the appropriate "fading memory" property when generated through a general dynamic binary choice model with regressors and possibly correlated errors. For the analysis of the smoothed maximum score estimator, this "fading memory" property that is proven for $y_{n}$ needs to be strong enough to allow a proof of an equivalent of the Hoeffding inequality, and in addition, it needs to allow for a proof of a central limit theorem (CLT) for a function of $y_{n}$ and $x_{n}$ that depends on $N$ in a situation where no martingale difference CLT can be applied. For a proof of validity of the dynamic probit model, the "fading memory" property only needs to support laws of large numbers and uniform laws of large numbers.
The "fading memory" property that we will prove for $y_{n}$ is that of near epoch dependence. The idea of the proof is similar to that of proofs for showing fading memory properties of

[^1]processes $y_{n}$ of the form
\[

$$
\begin{equation*}
y_{n}=f\left(y_{n-1}\right)+\varepsilon_{n}, \tag{5}
\end{equation*}
$$

\]

where $f($.$) is such that |f(x)-f(y)| \leq L|x-y|$ for some $L<1$. Functions $f($.$) satisfying this$ condition are called contraction mappings. Such proofs can be found in Bierens (1981) and Pötscher and Prucha (1997), for example. Pötscher and Prucha (1997, Section 6.4) contains a thorough discussion of these types of results, but the approach in the proof of this paper is different from the techniques discussed there. The differences are that the $f($.$) function$ in the dynamic binary choice case is not continuous, depends on $\varepsilon_{n}$, is not strictly less than 1 , and depends on more than one lagged value of $y_{n}$. These problems are essentially solved by smoothing the response function by the expectations operator, by using the fact that $y_{n}$ is a binary random variable, and by the use of the appropriate metric on the arguments of the $f($.$) function.$
Near epoch dependence of random variables $y_{n}$ on a base process of random variables $\eta_{n}$ is defined as follows:

Definition 1 Random variables $y_{n}$ are called near epoch dependent on $\eta_{n}$ if

$$
\begin{equation*}
\sup _{n \in \mathbb{Z}} E\left|y_{n}-E\left(y_{n} \mid \eta_{n-m}, \eta_{n-m+1}, \ldots, \eta_{n}\right)\right|^{2}=\nu(m)^{2} \rightarrow 0 \quad \text { as } \quad m \rightarrow \infty \tag{6}
\end{equation*}
$$

The base process $\eta_{n}$ needs to satisfy a condition such as strong or uniform mixing or independence. For the definitions of strong ( $\alpha-$ ) and uniform ( $\phi$-) mixing see e.g. Gallant and White (1988, p. 23) or Pötscher and Prucha (1997, p. 46). The near epoch dependence condition functions as a device that allows approximation of $y_{n}$ by a function of finitely many mixing or independent random variables $\eta_{n}$. Note also that for strictly stationary $\left(y_{n}, \eta_{n}\right)$, the "sup" in the above definition can be removed, because in that case

$$
\begin{equation*}
E\left|y_{n}-E\left(y_{n} \mid \eta_{n-m}, \eta_{n-m+1}, \ldots, \eta_{n}\right)\right|^{2} \tag{7}
\end{equation*}
$$

does not depend on $n$. The reader is referred to Gallant and White (1988) for a detailed account of the near epoch dependence condition. See also Pötscher and Prucha (1997) for a more up-to-date treatment of dependence concepts such as near epoch dependence.
The main results of the paper are the conditions under which $y_{n}$ is stationary and near epoch dependent (Theorem 1) and the conditions under which $y_{n}$ is strong mixing (Theorem 2). Unlike the linear model autoregressive model, no restrictions on the parameter space are needed for stationarity, near epoch dependence or strong mixing.
For establishing near epoch dependence of $y_{n}$, we have the following result. Define $S$ as the set of all $2^{p}$ possible $p$-vectors $s$ such that its elements $s_{i}$ are 0 or 1 , and define

$$
\begin{equation*}
\boldsymbol{\Phi}=\left\{\phi: \phi=\sum_{i=1}^{p} \rho_{i} s_{i}, s \in S\right\} . \tag{8}
\end{equation*}
$$

Let $\phi_{\min }$ denote the smallest element of $\boldsymbol{\Phi}$, and let $\phi_{\max }$ denote the largest element.
Theorem 1 Assume that $y_{n}$ is generated as $y_{n}=I\left(\sum_{j=1}^{p} \rho_{j} y_{n-j}+\eta_{n}>0\right)$. Let $\eta_{n}$ be strong mixing and strictly stationary. Assume that there is some $\delta>0$ for which there exists a positive integer $K$ such that

$$
\begin{align*}
& P\left(\phi_{\max }+\max _{i=1, \ldots, p} \eta_{n-i}>0 \mid y_{n-p-K}, y_{n-p-K-1}, \ldots\right) \\
& \quad-P\left(\phi_{\min }+\min _{i=1, \ldots, p} \eta_{n-i}>0 \mid y_{n-p-K}, y_{n-p-K-1}, \ldots\right)<1-\delta \quad \text { almost surely. } \tag{9}
\end{align*}
$$

Then (i) $y_{n}$ is near epoch dependent on $\eta_{n}$ and the near epoch dependence sequence $\nu($. satisfies $\nu(m) \leq C_{1} \exp \left(-C_{2} m\right)$, for positive constants $C_{1}$ and $C_{2}$; (ii) $\left(y_{n}, \eta_{n}\right)$ is strictly stationary.
Note that if $\eta_{n}=\gamma^{\prime} x_{n}+u_{n}$ for strong mixing and strictly stationary $\left(x_{n}^{\prime}, u_{n}\right)$, clearly $\eta_{n}$ is mixing as well. This observation will be used below. The assumption of Equation (9) limits the predictability of $y_{t}$ given the distant past. From the definition of uniform mixing, by choosing $K$ large enough, the assumption of Equation (9) follows if $\eta_{n}$ is uniform mixing and

$$
\begin{equation*}
P\left(\phi_{\max }+\max _{i=1, \ldots, p} \eta_{n-i}>0\right)-P\left(\phi_{\min }+\min _{i=1, \ldots, p} \eta_{n-i}>0\right)<1-\delta \tag{10}
\end{equation*}
$$

The proof of the above result is substantially easier for the case where $\eta_{n}$ is i.i.d., only one lagged $y_{n}$ is used as regressor and no other regressors are included. In that case, we can write

$$
\begin{equation*}
y_{n}=y_{n-1} I\left(\rho_{1}+\eta_{n}>0\right)+\left(1-y_{n-1}\right) I\left(\eta_{n}>0\right) \tag{11}
\end{equation*}
$$

implying that

$$
\begin{align*}
& \nu_{m} \equiv \sup _{n \in \mathbb{Z}} E\left|y_{n}-E\left(y_{n} \mid \eta_{n-m}, \ldots, \eta_{n}\right)\right|^{2} \\
& =\sup _{n \in \mathbb{Z}} E \mid\left(I\left(\rho_{1}+\eta_{n}>0\right)-I\left(\eta_{n}>0\right)\right)\left(y_{n-1}-\left.E\left(y_{n-1} \mid \eta_{n-m}, \ldots, \eta_{n-1}\right)\right|^{2}\right. \\
& =\left|P\left(\rho_{1}+\eta_{n}>0\right)-P\left(\eta_{n}>0\right)\right| \sup _{n \in \mathbb{Z}} E\left|y_{n-1}-E\left(y_{n-1} \mid \eta_{(n-1)-(m-1)}, \ldots, \eta_{n-1}\right)\right|^{2} \\
& =\left|P\left(\rho_{1}+\eta_{n}>0\right)-P\left(\eta_{n}>0\right)\right| \nu_{m-1} \tag{12}
\end{align*}
$$

which implies that the $\nu(m)$ sequence decays exponentially under the condition of Equation (9). The proof of Theorem 1 should be viewed as an extension to the above reasoning.

The fact that $y_{n}$ is a $0 / 1$-valued near epoch dependent random variable can now be exploited to show that $\left(y_{n}, x_{n}^{\prime}\right)$ is also strong mixing. Note that this is an observation that apparently has not been made in the literature before. The result is as follows:

Theorem 2 Suppose that $y_{n}=f\left(\eta_{n}, \eta_{n-1}, \ldots\right)$ is a sequence of 0/1-valued random variables that is near epoch dependent on $\left(u_{n}, x_{n}^{\prime}\right)$ with near epoch dependence coefficients $\nu(m)$, where $\eta_{n}=\gamma^{\prime} x_{n}+u_{n}$ and $\left(u_{n}, x_{n}^{\prime}\right)^{\prime}$ is strong mixing with mixing coefficients $\alpha(m)$. Then $\left(y_{n}, x_{n}^{\prime}\right)^{\prime}$ is strong mixing with strong mixing coefficients $C(\nu(m)+\alpha(m))$ for some $C>0$.

The mixing property of $\left(y_{n}, x_{n}^{\prime}\right)$ will be used in the proofs for consistency and asymptotic normality of the next sections.

## 3 The dynamic probit model

This section examines the behavior of the dynamic probit model estimator that results from including lagged $y_{n}$ among the regressors. Let $\beta=\left(\rho^{\prime}, \gamma^{\prime}\right)^{\prime}$ denote the true parameter and let $b=\left(r^{\prime}, c^{\prime}\right)^{\prime}, \rho, r \in \mathbb{R}^{p}$ and $\gamma, c \in \mathbb{R}^{q}$, and let $R$ and $\Gamma$ denote the parameter spaces for $r$ and $c$ respectively, and let $B=R \times \Gamma$. We assume normality of the errors so that the normalized loglikelihood conditional on $y_{1}, \ldots, y_{p}$ has the following form,

$$
\begin{align*}
& L_{N}(b)=(N-p)^{-1} \sum_{n=p+1}^{N} l_{n}(b) \\
= & (N-p)^{-1} \sum_{n=p+1}^{N}\left[y_{n} \log \left(\Phi\left(\sum_{j=1}^{p} r_{j} y_{n-j}+c^{\prime} x_{n}\right)\right)+\left(1-y_{n}\right) \log \left(1-\Phi\left(\sum_{j=1}^{p} r_{j} y_{n-j}+c^{\prime} x_{n}\right)\right)\right] . \tag{13}
\end{align*}
$$

Given the result of Theorem 2, it is now straightforward to find standard conditions under which the maximum likelihood estimator $b_{N}^{M L}$ is consistent.

## Assumption A

1. $x_{n}$ is a sequence of strictly stationary strong mixing random variables with $\alpha$-mixing numbers $\alpha(m)$, where $x_{n} \in \mathbb{R}^{q}$ for $q \geq 0$ and $\gamma \in \mathbb{R}^{q}$, and the second absolute moment of $x_{n}$ exits. The distribution of $w_{n}=\left(x_{n}^{\prime}, y_{n-1}, \ldots, y_{n-p}\right)^{\prime}$ is not contained in any linear subspace of $\mathbb{R}^{q}$.
2. $u_{n} \mid x_{n}, y_{n-1}, y_{n-2} \ldots, y_{n-p} \sim \operatorname{iid} N(0,1)$.
3. $y_{n}=I\left(\sum_{i=1}^{p} \rho_{i} y_{n-i}+\gamma^{\prime} x_{n}+u_{n}>0\right)$.
4. $\beta$ is an element of the interior of a convex set $B$.

Theorem 3 Under Assumption $A, b_{N}^{M L} \xrightarrow{p} \beta$. If in addition (i) the strong mixing coefflcients satisfy $\alpha(m) \leq C m^{-\eta}$ for positive constants $C$ and $\eta$ and (ii) $E\left|l_{n}(b)\right|^{1+\delta}<\infty$ for some $\delta>0$ and all $b \in B$, and (iii) $B$ is compact, then $b_{N}^{M L} \xrightarrow{\text { as }} \beta$.

Let $I=-E(\partial / \partial b)\left(\partial / \partial b^{\prime}\right) l_{n}(\beta)$. For asymptotic normality, we need an additional assumption.

## Assumption B

1. $u_{n} \mid\left(x_{n}, y_{n-1}\right),\left(x_{n-1}, y_{n-2}\right) \ldots \sim \operatorname{iid} N(0,1)$.

Theorem 4 Under Assumptions $A$ and $B, N^{1 / 2}\left(b_{N}^{M L}-\beta\right) \xrightarrow{d} N\left(0, I^{-1}\right)$.
Under the above Assumptions A and B, it also follows that the usual estimators of $I$, using either the outer product or Hessian approach, will both be weakly consistent for $I$.
Note that given the weak dependence property of Theorem 2, it is also possible to set forth conditions such that for weakly dependent $u_{n}$ with arbitrary distribution, $N^{1 / 2}\left(b_{N}^{M L}-\beta^{*}\right) \xrightarrow{d}$ $N(0, J)$ for some matrix $J$ and a $\beta^{*}$ that uniquely minimizes the objective function. Here of course $\beta^{*}$ does not necessarily equal the true parameter value $\beta$. However, in order to show that the probit objective function is uniquely maximized at $\beta$, we need that a first order condition of the type

$$
\begin{equation*}
E\left(y_{n}-\Phi\left(\sum_{i=1}^{p} \rho_{i} y_{n-i}+\gamma^{\prime} x_{n}\right)\right) m\left(y_{n-1}, \ldots, y_{n-p}, x_{n}\right)=0 \tag{14}
\end{equation*}
$$

holds for some function $m(., \ldots,$.$) . This condition is implied by$

$$
\begin{equation*}
E\left(y_{n} \mid y_{n-1}, \ldots, x_{n}\right)=\Phi\left(\sum_{i=1}^{p} \rho_{i} y_{n-i}+\gamma^{\prime} x_{n}\right) \tag{15}
\end{equation*}
$$

and the latter condition is equivalent to assuming that $u_{n}$ is i.i.d. and standard normal if lagged values of $y_{n}$ are included.

## 4 Consistency of the smoothed maximum score estimator

The smoothed maximum score estimator is defined as $\operatorname{argmax}_{b \in B} S_{N}\left(b, \sigma_{N}\right)$, where

$$
\begin{equation*}
S_{N}\left(b, \sigma_{N}\right)=(N-p)^{-1} \sum_{n=p+1}^{N}\left(2 \cdot I\left(y_{n}=1\right)-1\right) K\left(\left(\sum_{j=1}^{p} r_{j} y_{n-j}+c^{\prime} x_{n}\right) / \sigma_{N}\right) \tag{16}
\end{equation*}
$$

and $\sigma_{N}$ is a bandwidth-type sequence such that $\sigma_{N} \rightarrow 0$ as $N \rightarrow \infty$, where $K($.$) is a function$ such that $K(-\infty)=0$ and $K(\infty)=1$. This objective function is a smoothed version of the maximum score objective function

$$
\begin{equation*}
S_{N}^{*}(b)=(N-p)^{-1} \sum_{n=p+1}^{N}\left(2 \cdot I\left(y_{n}=1\right)-1\right) I\left(\sum_{j=1}^{p} r_{j} y_{n-j}+c^{\prime} x_{n} \geq 0\right) \tag{17}
\end{equation*}
$$

In addition, let $S(b)=E S_{N}^{*}(b)$. This notation is justified because we will use conditions under which $\left(y_{n}, x_{n}\right)$ will be proven to be strictly stationary. See Manski (1985) and Kim and Pollard (1990) for more information and results regarding the maximum score estimator. While Horowitz' maximum score estimator can reach the optimal rate of convergence (see Horowitz (1992)), Kim and Pollard (1990) showed that the maximum score estimator in general is consistent of order $N^{-1 / 3}$.
The following five assumptions are needed for the proof of our consistency result:
Assumption $1\left(x_{n}^{\prime}, u_{n}\right)^{\prime}$ is a sequence of strictly stationary strong mixing random variables with $\alpha$-mixing numbers $\alpha(m)$, where $x_{n} \in \mathbb{R}^{q}$ for $q \geq 1$ and $\gamma \in \mathbb{R}^{q}$, and

$$
\begin{equation*}
y_{n}=I\left(\sum_{i=1}^{p} \rho_{i} y_{n-i}+\gamma^{\prime} x_{n}+u_{n}>0\right) . \tag{18}
\end{equation*}
$$

Note that by Theorem 1 and the discussion following that theorem, $\left(y_{n}, x_{n}\right)^{\prime}$ is strictly stationary. This justifies the formulation of the assumptions below in their current forms. Define $\tilde{x}_{n}=\left(y_{n-1}, \ldots, y_{n-p}, x_{n 2}, \ldots, x_{n q}\right)$.

Assumption 2 The support of the distribution of $\left(x_{n 1}, \tilde{x}_{n}^{\prime}\right)^{\prime}$ is not contained in any proper linear subspace of $\mathbb{R}^{p+q}$. (b) $0<P\left(y_{n}=1 \mid x_{n 1}, \tilde{x}_{n}\right)<1$ almost surely. (c) $\gamma_{1} \neq 0$, and for almost every $\tilde{x}_{n}$, the distribution of $x_{n 1}$ conditional on $\tilde{x}_{n}$ has everywhere positive density with respect to Lebesgue measure.

Assumption $3 \operatorname{Median}\left(u_{n} \mid x_{n}, y_{n-1}, \ldots, y_{n-p}\right)=0 \quad$ almost surely.
Assumption 3 allows for heteroskedasticity of arbitrary form, including heteroskedasticity that depends on lagged values of $y_{n}$. If all regressors are exogenous, Assumption 3 allows for correlated errors, e.g. the errors could follow an ARMA process.

Assumption $4\left|\gamma_{1}\right|=1$, and $\tilde{\beta}=\left(\rho_{1}, \ldots, \rho_{p}, \gamma_{2}, \ldots, \gamma_{q}\right)$ is contained in a compact subset $\tilde{B}$ of $\mathbb{R}^{p+q-1}$.

The following assumption is simply the assumption of Equation (9) in Theorem 1 for $\eta_{n}=$ $\gamma^{\prime} x_{n}+u_{n}$.

Assumption 5 For $\phi_{\max }$ and $\phi_{\min }$ as defined before, for some $\delta>0$ there exists a positive integer $K$ such that

$$
\begin{align*}
& P\left(\phi_{\max }+\max _{i=1, \ldots, p}\left(\gamma^{\prime} x_{n-i}+u_{n-i}\right)>0 \mid y_{n-p-K}, y_{n-p-K-1}, \ldots\right) \\
& \quad-P\left(\phi_{\min }+\min _{i=1, \ldots, p}\left(\gamma^{\prime} x_{n-i}+u_{n-i}\right)>0 \mid y_{n-p-K}, y_{n-p-K-1}, \ldots\right)<1-\delta . \tag{19}
\end{align*}
$$

We need some form of scale normalization; we set $\left|b_{1}\right|=1$ here, as in Horowitz (1992). Therefore, the estimator $b_{N}$ needs to be defined as

$$
\begin{equation*}
b_{N}=\operatorname{argmax}_{b:\left|b_{1}\right|=1} S_{N}\left(b, \sigma_{N}\right) \tag{20}
\end{equation*}
$$

The following result shows the consistency of $b_{N}$ :
Theorem 5 Under Assumptions 1,2,3,4 and $5, b_{N} \xrightarrow{p} \beta$. If in addition the strong mixing coefficients satisfy $\alpha(m) \leq C m^{-\eta}$ for positive constants $C$ and $\eta$, then $b_{N} \xrightarrow{\text { as }} \beta$.

## 5 Asymptotic normality of the smoothed maximum score estimator

Define, analogously to Horowitz (1992), $\tilde{b}=\left(r_{1}, \ldots, r_{p}, c_{2}, \ldots, c_{q}\right)$, and let

$$
\begin{align*}
& T_{N}\left(b, \sigma_{N}\right)=\partial S_{N}\left(b, \sigma_{N}\right) / \partial \tilde{b}  \tag{21}\\
& Q_{N}\left(b, \sigma_{N}\right)=\partial^{2} S_{N}\left(b, \sigma_{N}\right) / \partial \tilde{b} \partial \tilde{b}^{\prime} \tag{22}
\end{align*}
$$

Also, define

$$
\begin{equation*}
z_{n}=\sum_{j=1}^{p} \rho_{j} y_{n-j}+\gamma^{\prime} x_{n} \tag{23}
\end{equation*}
$$

and let $p\left(z_{n} \mid \tilde{x}_{n}\right)$ denote the density of $z_{n}$ given $\tilde{x}_{n}$, let $P($.$) denote the distribution of \tilde{x}_{n}$, let $F\left(. \mid z_{n}, \tilde{x}_{n}\right)$ denote the cumulative distribution of $u_{n}$ conditional on $z_{n}$ and $\tilde{x}_{n}$. For each positive integer $i$, define

$$
\begin{equation*}
F^{(i)}(-z, x, \tilde{x})=\partial^{i} F(-z \mid z, \tilde{x}) / \partial z^{i} \tag{24}
\end{equation*}
$$

Let $h$ denote a positive integer that satisfies the conditions of Assumptions 8, 9 and 10 below, and let

$$
\begin{align*}
& \alpha_{A}=\int_{-\infty}^{\infty} v^{h} K^{\prime}(v) d v  \tag{25}\\
& \alpha_{D}=\int_{-\infty}^{\infty} K^{\prime}(v)^{2} d v \tag{26}
\end{align*}
$$

Also analogously to Horowitz (1992), define

$$
\begin{align*}
& A=-2 \alpha_{A} \sum_{i=1}^{h}\left\{[i!(h-i)!]^{-1} E\left[F^{(i)}\left(0,0, \tilde{x}_{n}\right) p^{(h-i)}\left(0 \mid \tilde{x}_{n}\right) \tilde{x}_{n}\right]\right\},  \tag{27}\\
& D=\alpha_{D} \cdot E\left[\tilde{x}_{n} \tilde{x}_{n}^{\prime} p\left(0 \mid \tilde{x}_{n}\right)\right]  \tag{28}\\
& Q=2 \cdot E\left[\tilde{x}_{n} \tilde{x}_{n}^{\prime} F^{(1)}\left(0 \mid 0, \tilde{x}_{n}\right) p\left(0 \mid \tilde{x}_{n}\right)\right] . \tag{29}
\end{align*}
$$

The following assumption is the analogue of Horowitz' Assumption 5, which is the assumption below for $s=4$. It appears that Horowitz' truncation argument is in error (see also notes 2, 3,4 and 5), but that his argument is correct for bounded data. This explains the presence here of a condition that is stronger than that of Horowitz.

Assumption 6 For all vectors $\xi$ such that $|\xi|=1, E\left|\xi^{\prime} \tilde{x}\right|^{s}<\infty$ for some $s>4$.
We need to strengthen the fading memory conditions of Assumption 1 in order to establish asymptotic normality:

Assumption 1' $\left(x_{n}^{\prime}, u_{n}\right)$ is a sequence of strictly stationary strong mixing random variables with $\alpha$-mixing numbers $\alpha(m)$ such that $\alpha(m) \leq C m^{-(2 s-2) /(s-2)-\eta}$ for some $\eta>0$, where $x_{n} \in \mathbb{R}^{q}$ for $q \geq 1$ and $\gamma \in \mathbb{R}^{q}$, and

$$
\begin{equation*}
y_{n}=I\left(\sum_{i=1}^{p} \rho_{i} y_{n-i}+\gamma^{\prime} x_{n}+u_{n}>0\right) . \tag{30}
\end{equation*}
$$

The assumption below is needed in lieu of Horowitz' Assumption 6.
Assumption 7 For some sequence $m_{N} \geq 1$,

$$
\begin{align*}
& \sigma_{N}^{-3(p+q-1)} \sigma_{N}^{-2} N^{1 / s} \alpha\left(m_{N}\right)+\sigma_{N}^{-2(p+q-1) / \beta} N^{2 / s} \alpha\left(m_{N}\right) \\
& \quad+\left|\log \left(N m_{N}\right)\right|\left(N^{1-4 / s} \sigma_{N}^{4} m_{N}^{-2}\right)^{-1} \rightarrow 0 \quad \text { as } \quad N \rightarrow \infty . \tag{31}
\end{align*}
$$

For the case of independent $\left(x_{n}, u_{n}\right), \alpha(m)=0$ for $m \geq 1$, and we can set $m_{N}=1$ for that case. The condition of Assumption 7 then becomes

$$
\begin{equation*}
(\log (N))\left(N^{1-4 / s} \sigma_{N}^{4}\right)^{-1} \rightarrow 0 \quad \text { as } \quad N \rightarrow \infty, \tag{32}
\end{equation*}
$$

implying that for bounded data, we can set $s=\infty$ and obtain Horowitz' condition

$$
\begin{equation*}
(\log (N))\left(N \sigma_{N}^{4}\right)^{-1} \rightarrow 0 \quad \text { as } \quad N \rightarrow \infty \tag{33}
\end{equation*}
$$

The following assumptions are identical to Horowitz' Assumptions 7-11:
Assumption 8 (a) $K($.$) is twice differentiable everywhere, |K()$.$| and K^{\prime \prime}($.$) are uniformly$ bounded, and each of the following integrals over $(-\infty, \infty)$ is finite: $\int\left[K^{\prime}(v)\right]^{4} d v, \int\left[K^{\prime \prime}(v)\right]^{2} d v$, $\int\left|v^{2} K^{\prime \prime}(v)\right| d v$. (b) For some integer $h \geq 2$ and each integer $i(1 \leq i \leq h), \int\left|v^{i} K^{\prime}(v)\right| d v<$ $\infty$, and

$$
\int_{-\infty}^{\infty} v^{i} K^{\prime}(v) d v= \begin{cases}0 & \text { if } i<h,  \tag{34}\\ d & \text { (nonzero) if } i=h .\end{cases}
$$

(c) For any integer $i$ between 0 and $h$, any $\eta>0$, and any sequence $\left\{\sigma_{N}\right\}$ converging to 0 ,

$$
\begin{equation*}
\lim _{N \rightarrow \infty} \sigma_{N}^{i-h} \int_{\left|\sigma_{N} v\right|>\eta}\left|v^{i} K^{\prime}(v)\right| d v=0 \tag{35}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{N \rightarrow \infty} \sigma_{N} \int_{\left|\sigma_{N} v\right|>\eta}\left|K^{\prime \prime}(v)\right| d v=0 \tag{36}
\end{equation*}
$$

Assumption 9 For each integer $i$ such that $1 \leq i \leq h-1$, all $z$ in a neighborhood of 0 , almost every $\tilde{x}_{n}$, and some $M<\infty, p^{(i)}\left(z_{n} \mid \tilde{x}_{n}\right)$ exists and is a continuous function of $z_{n}$ satisfying $\left|p^{(i)}\left(z_{n} \mid \tilde{x}_{n}\right)\right|<M$. In addition, $p\left(z_{n} \mid \tilde{x}_{n}\right)<M$ for all $z$ and almost every $\tilde{x}$.

Assumption 10 For each integer $i$ such that $1 \leq i \leq h$, all $z_{n}$ in a neighborhood of 0 , almost every $\tilde{x}_{n}$, and some $M<\infty, F^{(i)}\left(-z_{n}, z_{n}, \tilde{x}_{n}\right)$ exists and is a continuous function of $z_{n}$ satisfying $\left|F^{(i)}\left(-z_{n}, z_{n}, \tilde{x}_{n}\right)\right|<M$.

Assumption $11 \tilde{\beta}$ is an interior point of $\tilde{B}$.
Assumption 12 The matrix $Q$ is negative definite.
In addition to the above equivalents to Horowitz' assumptions, we will also need the following two assumptions. The first assumption is needed to assure proper behavior of covariance terms.

Assumption 13 The conditional joint density $p\left(z_{n}, z_{n-j} \mid x_{n}, x_{n-j}\right)$ exists for all $j \geq 1$ and is continuous at $\left(z_{n}, z_{n-j}\right)=(0,0)$ for all $j \geq 1$.

The next condition on $K^{\prime \prime}($.$) is needed to formally show a uniform law of large numbers for$ the second derivative of the objective function.

Assumption $14 K^{\prime \prime}($.$) satisfies, for some \mu \in(0,1]$ and $L \in[0, \infty)$ and all $x, y \in \mathbb{R}$,

$$
\begin{equation*}
\left|K^{\prime \prime}(x)-K^{\prime \prime}(y)\right| \leq L|x-y|^{\mu} . \tag{37}
\end{equation*}
$$

To prove asymptotic normality, we need an inequality in the spirit of Hoeffding's inequality, but for weakly dependent random variables. We derive such an inequality in the Appendix as Lemma 10. The inequality of Lemma 10 also allows for martingale difference sequences so that it covers both the random sample case of Horowitz (1992) as well as the dynamic case.
Our asymptotic normality result now is the following. This result, of course, is nearly identical to Horowitz' in the non-dynamic cross-section case.

Theorem 6 Let Assumptions 1' and Assumptions 2-14 hold for some $h \geq 2$. Then

1. If $N \sigma_{N}^{2 h+1} \rightarrow \infty$ as $N \rightarrow \infty, \sigma_{N}^{-h}\left(\tilde{b}_{N}-\tilde{\beta}\right) \xrightarrow{p}-Q^{-1} A$.
2. If $N \sigma_{N}^{2 h+1}$ has a finite limit $\lambda$ as $N \rightarrow \infty$,

$$
\begin{equation*}
\left(N \sigma_{N}\right)^{1 / 2}\left(\tilde{b}_{N}-\tilde{\beta}\right) \xrightarrow{d} N\left(-\lambda^{1 / 2} Q^{-1} A, Q^{-1} D Q^{-1}\right) . \tag{38}
\end{equation*}
$$

In order to estimate the matrices $A, D$ and $Q$, we need an additional result, the analogue of Horowitz' (1992) Theorem 3.

Theorem 7 Let $b_{N}$ be a consistent smoothed maximum score estimator based on $\sigma_{N}$ such that $\sigma_{N}=O\left(n^{-1 /(2 h+1)}\right)$. For $b \in\{-1,1\} \times \tilde{B}$, define

$$
\begin{equation*}
t_{n}(b, \sigma)=\left(2 \cdot I\left(y_{n}=1\right)-1\right)\left(\tilde{x}_{n} / \sigma_{N}\right) K^{\prime}\left(\left(\sum_{j=1}^{p} r_{j} y_{n-j}+c^{\prime} x_{n}\right) / \sigma\right) . \tag{39}
\end{equation*}
$$

Let $\sigma_{N}^{*}$ be such that $\sigma_{N}^{*}=O\left(N^{-\delta /(2 h+1)}\right)$, where $0<\delta<1$. Then: (a) $\hat{A}_{N} \equiv\left(\sigma_{N}^{*}\right)^{-h} T_{N}\left(b_{N}, \sigma_{N}^{*}\right)$ converges in probability to $A$; (b) the matrix

$$
\begin{equation*}
\hat{D}_{N} \equiv \sigma_{N}(N-p)^{-1} \sum_{n=p+1}^{N} t_{n}\left(b_{N}, \sigma_{N}\right) t_{n}\left(b_{N}, \sigma_{N}\right)^{\prime} \tag{40}
\end{equation*}
$$

converges in probability to $D$; (c) $Q_{N}\left(b_{N}, \sigma_{N}\right)$ converges in probability to $Q$.

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## Proofs

## Proof of Theorem 1:

The dynamic binary choice model of order $p$ can be written as

$$
y_{n}=I\left(\sum_{i=1}^{p} \rho_{i} y_{n-i}+\eta_{n}>0\right)=g\left(y_{n-1}, y_{n-2}, \ldots, y_{n-p}, \eta_{n}\right) .
$$

This $g(., \ldots,$.$) satisfies, for all 0-1 valued y_{1}, y_{2}, \ldots, y_{n-p}$ and $\tilde{y}_{1}, \tilde{y}_{2}, \ldots, \tilde{y}_{n-p}$,

$$
\begin{aligned}
& \left|g\left(y_{1}, y_{2}, \ldots, y_{p}, \eta_{n}\right)-g\left(\tilde{y}_{1}, \tilde{y}_{2}, \ldots, \tilde{y}_{p}, \eta_{n}\right)\right| \leq L\left(\eta_{n}\right) \max _{j=1, \ldots, p}\left|y_{j}-\tilde{y}_{j}\right|, \\
& L\left(\eta_{n}\right)=\sup _{\phi, \phi^{\prime} \in \boldsymbol{\Phi}}\left|I\left(\phi+\eta_{n}>0\right)-I\left(\phi^{\prime}+\eta_{n}>0\right)\right| .
\end{aligned}
$$

where $\boldsymbol{\Phi}$ was defined in Equation (8). The idea of the proof is to show that the process $y_{n}$ can be approximated arbitrarily well by using a function of a finite number of $\eta_{n}$ - this is the content of the near epoch dependence concept. We do this by using for our approximation $\hat{y}_{n}^{m}$ the $y$ that would have resulted if the process had been started up using 0 values for the $y_{n}$ and $\eta_{n}$ that occurred $m$ periods or longer ago. Formally, for all $n$ define $\hat{y}_{n}^{m}=0$ for $m \leq 0$. Then for all $m \geq 1$ recursively define

$$
\hat{y}_{n}^{m}=g\left(\hat{y}_{n-1}^{m-1}, \hat{y}_{n-2}^{m-2}, \ldots, \hat{y}_{n-p}^{m-p}, \eta_{n}\right) .
$$

Note that by construction, $\hat{y}_{n}^{m}=f_{m}\left(\eta_{n}, \eta_{n-1}, \ldots, \eta_{n-m}\right)$. Define $\max _{j \in A} c_{j}=0$ if $A$ is empty. Then for these approximators $\hat{y}_{n}^{m}$ we have, using $0 \leq L(\cdot) \leq 1$,

$$
\begin{aligned}
& \max _{j=1, \ldots, p}\left|y_{n-j+1}-\hat{y}_{n-j+1}^{m-j+1}\right| \\
= & \max \left(\left|g\left(y_{n-1}, y_{n-2}, \ldots, y_{n-p}, \eta_{n}\right)-g\left(\hat{y}_{n-1}^{m-1}, \hat{y}_{n-2}^{m-2}, \ldots, \hat{y}_{n-p}^{m-p}, \eta_{n}\right)\right|, \max _{j=2, \ldots, p}\left|y_{n-j+1}-\hat{y}_{n-j+1}^{m-j+1}\right|\right) \\
& \leq \max \left(L\left(\eta_{n}\right) \max _{j=1, \ldots, p}\left|y_{n-j}-\hat{y}_{n-j}^{m-j}\right|, \max _{j=2, \ldots, p}\left|y_{n-j+1}-\hat{y}_{n-j+1}^{m-j+1}\right|\right) \\
& \leq \max \left(L\left(\eta_{n}\right)\left|y_{n-p}-\hat{y}_{n-p}^{m-p}\right|, \max _{j=2, \ldots, p}\left|y_{n-j+1}-\hat{y}_{n-j+1}^{m-j+1}\right|\right) \\
& \leq \max \left(L\left(\eta_{n}\right)\left|y_{n-p}-\hat{y}_{n-p}^{m-p}\right|, L\left(\eta_{n-1}\right) \max _{j=1, \ldots, p}\left|y_{n-j-1}-\hat{y}_{n-j-1}^{m-j-1}\right|, \max _{j=3, \ldots, p}\left|y_{n-j+1}-\hat{y}_{n-j+1}^{m-j+1}\right|\right) \\
& \leq \max \left(L\left(\eta_{n}\right)\left|y_{n-p}-\hat{y}_{n-p}^{m-p}\right|, L\left(\eta_{n-1}\right)\left|y_{n-p-1}-\hat{y}_{n-p-1}^{m-p-1}\right|,\right.
\end{aligned}
$$

$$
\begin{aligned}
& \left.L\left(\eta_{n-1}\right)\left|y_{n-p}-\hat{y}_{n-p}^{m-p}\right|, \max _{j=3, \ldots, p}\left|y_{n-j+1}-\hat{y}_{n-j+1}^{m-j+1}\right|\right) \\
\leq & \max _{j=1, \ldots, p} L\left(\eta_{n-j+1}\right) \max _{j=1, \ldots, p}\left|y_{n-p-j+1}-\hat{y}_{n-p-j+1}^{m-p-j+1}\right|
\end{aligned}
$$

and again using $0 \leq L(\cdot) \leq 1$, we also have by repeating this reasoning $K$ times, for all $K \geq 1$,

$$
\begin{aligned}
& \max _{j=1, \ldots, p}\left|y_{n-j+1}-\hat{y}_{n-j+1}^{m-j+1}\right| \\
& \leq \max _{j=1, \ldots, p} L\left(\eta_{n-j+1}\right) \max _{j=1, \ldots, p}\left|y_{n-p K-j+1}-\hat{y}_{n-p K-j+1}^{m-p K-j+1}\right|
\end{aligned}
$$

Next, note that by assumption there exists a positive integer $\tilde{K}$ such that, for some $\delta>0$, for $\phi_{\max }$ and $\phi_{\min }$ as defined below Equation (8),

$$
\begin{aligned}
& \left|E\left(\max _{j=1, \ldots, p} L\left(\eta_{n-j+1}\right) \mid y_{n-p \tilde{K}}, y_{n-p \tilde{K}-1}, \ldots\right)\right| \\
& =\mid P\left(\phi_{\max }+\max _{j=1, \ldots, p} \eta_{n-j+1}>0 \mid y_{n-p \tilde{K}}, y_{n-p \tilde{K}-1}, \ldots\right) \\
& \quad-P\left(\phi_{\min }+\min _{j=1, \ldots, p} \eta_{n-j+1}>0 \mid y_{n-p \tilde{K}}, y_{n-p \tilde{K}-1}, \ldots\right) \mid<1-\delta
\end{aligned}
$$

for some $\delta>0$. Therefore, for $m \geq p \tilde{K}$,

$$
\begin{aligned}
& \chi_{m} \stackrel{\text { def }}{=} \sup _{n \in \mathbb{Z}} E \max _{j=1, \ldots, p}\left|y_{n-j+1}-\hat{y}_{n-j+1}^{m-j+1}\right|^{2} \\
& \leq \sup _{n \in \mathbb{Z}} E\left(E\left(\max _{j=1, \ldots, p} L\left(\eta_{n-j+1}\right) \mid y_{n-p \tilde{K}}, y_{n-p \tilde{K}-1}, \ldots\right) \max _{j=1, \ldots, p}\left|y_{n-p \tilde{K}-j+1}-\hat{y}_{n-p \tilde{K}-j+1}^{m-p \tilde{K}-j+1}\right|^{2}\right) \\
& \leq(1-\delta) \sup _{n \in \mathbb{Z}} E \max _{j=1, \ldots, p}\left|y_{n-p \tilde{K}-j+1}-\hat{y}_{n-p \tilde{K}-j+1}^{m-\tilde{\tilde{K}}-j+1}\right|^{2} \\
& =(1-\delta) \sup _{n \in \mathbb{Z}} E \max _{j=1, \ldots, p}\left|y_{n-j+1}-\hat{y}_{n-j+1}^{m-p \tilde{K}-j+1}\right|^{2} \\
& =(1-\delta) \chi_{m-p \tilde{K}} .
\end{aligned}
$$

The one but last equality follows because $\sup _{n \in \mathbb{Z}} f(n-p \tilde{K})=\sup _{n \in \mathbb{Z}} f(n)$ for any function $f(\cdot)$. Because $\chi_{j} \leq 1$ for all $j \geq 0$, it now follows that

$$
\chi_{m} \leq(1-\delta)^{[m /(p \tilde{K})]}
$$

where $[x]$ denotes the integer part of $x$. Now, because the conditional expectation is the best possible $L_{2}$-approximation,

$$
\begin{aligned}
& \sup _{n \in \mathbb{Z}} E\left|y_{n}-E\left(y_{n} \mid \eta_{n}, \eta_{n-1}, \ldots, \eta_{n-m}\right)\right|^{2} \\
& \leq \sup _{n \in \mathbb{Z}} E \max _{j=1, \ldots, p}\left|y_{n-j+1}-\hat{y}_{n-j+1}^{m-j+1}\right|^{2} \\
& \leq(1-\delta)^{[m /(2 \tilde{K})]} \leq C_{1} \exp \left(-C_{2} m\right)
\end{aligned}
$$

for positive constants $C_{1}$ and $C_{2}$.
To show that $\left(y_{n}, \eta_{n}\right)$ is strictly stationary, note that $\hat{y}_{n}^{m}=f_{m}\left(\eta_{n-m}, \ldots, \eta_{n}\right)$ by construction, where $f_{m}(., \ldots,$.$) does not depend on n$ or $N$. This then implies that $\left(y_{n}, \eta_{n}\right)$ is strictly stationary.

## Proof of Theorem 2:

Let $\mathcal{X}_{a, b}$ denote the $\sigma$-algebra generated by $\left(\left(x_{a}, y_{a}\right), \ldots,\left(x_{b}, y_{b}\right)\right)$. The definition of the strong mixing coefficients is
$\sup _{n \in \mathbb{Z}} \sup _{F \in \mathcal{X}_{-\infty, t}, G \in \mathcal{X}_{t+m, \infty}}\left\{\left|P\left(\left(x_{n}, y_{n}\right) \in F,\left(x_{n+m}, y_{n+m}\right) \in G\right)-P\left(\left(x_{n}, y_{n}\right) \in F\right) P\left(\left(x_{n+m}, y_{n+m}\right) \in G\right)\right|\right\}$,
see for example White (2001, page 47). Because $y_{n}$ is a $0 / 1$-valued random variable, there are only four possibilities for the possible values of the $\left(y_{n}, y_{n-m}\right)$ pair. Therefore,

$$
\begin{aligned}
& P\left(\left(x_{n}, y_{n}\right) \in F,\left(x_{n+m}, y_{n+m}\right) \in G\right) \\
& =E \sum_{i=0}^{1} \sum_{j=0}^{1} I\left(\left(x_{n}, y_{n}\right) \in F,\left(x_{n+m}, y_{n+m}\right) \in G\right) I\left(y_{n}=i\right) I\left(y_{n+m}=j\right) \\
& =E \sum_{i=0}^{1} \sum_{j=0}^{1} I\left(\left(x_{n}, i\right) \in F,\left(x_{n+m}, j\right) \in G\right) I\left(y_{n}=i\right) I\left(y_{n+m}=j\right)
\end{aligned}
$$

and

$$
P\left(\left(x_{n}, y_{n}\right) \in F\right)=E I\left(\left(x_{n}, y_{n}\right) \in F\right) \sum_{i=0}^{1} I\left(y_{n}=i\right)
$$

implying that

$$
\begin{aligned}
& \left|P\left(\left(x_{n}, y_{n}\right) \in F,\left(x_{n+m}, y_{n+m}\right) \in G\right)-P\left(\left(x_{n}, y_{n}\right) \in F\right) P\left(\left(x_{n+m}, y_{n+m}\right) \in G\right)\right| \\
& \leq \sum_{i=0}^{1} \sum_{j=0}^{1} \mid E I\left(\left(x_{n}, i\right) \in F,\left(x_{n+m}, j\right) \in G\right) I\left(y_{n}=i\right) I\left(y_{n+m}=j\right) \\
& \quad-E I\left(\left(x_{n}, i\right) \in F\right) I\left(y_{n}=i\right) E I\left(\left(x_{n+m}, j\right) \in G\right) I\left(y_{n+m}=j\right) \mid
\end{aligned}
$$

For the case $y_{n}=1, y_{n+m}=1$, we now have, defining $\mathcal{F}_{n}=\sigma\left(v_{n}, v_{n-1}, \ldots\right)$ for $v_{n}=\left(u_{n}, x_{n}^{\prime}\right)^{\prime}$,

$$
\begin{aligned}
& \left|E I\left(\left(x_{n}, 1\right) \in F\right) y_{n} I\left(\left(x_{n+m}, 1\right) \in G\right) y_{n+m}-E I\left(\left(x_{n}, 1\right) \in F\right) y_{n} E I\left(\left(x_{n+m}, 1\right) \in G\right) y_{n+m}\right| \\
& =\left|E I\left(\left(x_{n}, 1\right) \in F\right) y_{n}\left[E\left(I\left(\left(x_{n+m}, 1\right) \in G\right) y_{n+m} \mid \mathcal{F}_{n}\right)-E I\left(\left(x_{n+m}, 1\right) \in G\right) y_{n+m}\right]\right| \\
& \leq E\left|E\left(\left(I\left(x_{n+m}, 1\right) \in G\right) y_{n+m} \mid \mathcal{F}_{n}\right)-E\left(I\left(\left(x_{n+m}, 1\right) \in G\right) y_{n+m}\right)\right|,
\end{aligned}
$$

and convergence to zero with $m$ of the last expression constitutes the $L_{1}$-mixingale condition for $I\left(\left(x_{n}, 1\right) \in G\right) y_{n}$ with respect to $\mathcal{F}_{n}$; see for example Pötscher and Prucha (1997) for a definition on an $L_{1}$-mixingale. Now $I\left(\left(x_{n}, 1\right) \in G\right) y_{n}$ is a sequence that is bounded and near epoch dependent on $v_{n}$, implying that it is an $L_{1}$-mixingale, which in turn implies that

$$
\begin{aligned}
& E\left|E\left(I\left(\left(x_{n}, 1\right) \in G\right) y_{n} \mid \mathcal{F}_{n-m}\right)-E\left(I\left(\left(x_{n}, 1\right) \in G\right) y_{n}\right)\right| \\
& \leq C(\nu(m)+\alpha(m)) .
\end{aligned}
$$

The cases $y_{n}=1, y_{n+m}=0 ; y_{n}=0, y_{n+m}=1$; and $y_{n}=0, y_{n-m}=0$ are analogous, which then proves the result.

For the proof of Theorem 3, we need the following two lemmas. Let $w_{n}=\left(y_{n-1}, \ldots, y_{n-p}, x_{n}^{\prime}\right)^{\prime}$.
Lemma 1 Under the conditions of Theorem 3, and B being compact, $E \sup _{b \in B}\left|l_{n}(b)\right|<\infty$.

## Proof of Lemma 1:

Note that $E w_{n} w_{n}^{\prime}$ exists by Assumption A1. Existence of $E w_{n} w_{n}^{\prime}$ and the probit specification imply the result. The reasoning is similar to the result for cross-section probit, see Newey and McFadden (1994, page 2125, Example 1.2).

Lemma 2 Under the conditions of Theorem 3, (i) $E w_{n} w_{n}^{\prime}$ is positive definite and (ii) $E l_{n}(b)$ is uniquely maximized at $b=\beta$.

## Proof of Lemma 2:

Note that $E w_{n} w_{n}^{\prime}$ exists by Assumption A1. The assumptions of Theorem 1 are satisfied so that $\left(x_{n}^{\prime}, y_{n}\right)^{\prime}$ is strongly stationary. The assumption that distribution of $w_{n}$ is not contained in any linear subspace of $\mathbb{R}^{p+q}$ implies that $E w_{n} w_{n}^{\prime}$ is nonsingular so that $E w_{n} w_{n}^{\prime}$ is positive definite. Let $b \neq \beta$ so that $E\left[\left(w_{n}^{\prime}(b-\beta)\right)^{2}\right]=(b-\beta)^{\prime} E w_{n} w_{n}^{\prime}(b-\beta)>0$, implying that $w_{n}^{\prime}(b-\beta) \neq 0$ on a set with positive probability, implying that $w_{n}^{\prime} b \neq w_{n}^{\prime} \beta$ on a set with positive probability. Both $\Phi(z)$ and $\bar{\Phi}(z)=1-\Phi(z)$ are strictly monotonic, and therefore $w_{n}^{\prime} b \neq w_{n}^{\prime} \beta$ implies that both $\Phi\left(w_{n}^{\prime} b\right) \neq \Phi\left(w_{n}^{\prime} \beta\right)$ and $\bar{\Phi}\left(w_{n}^{\prime} b\right) \neq \Phi\left(w_{n}^{\prime} \beta\right)$. Thus, the density

$$
p\left(y_{n} \mid w_{n}, b\right)=\Phi\left(w_{n}^{\prime} b\right)^{y_{n}} \bar{\Phi}\left(w_{n}^{\prime} b\right)^{1-y_{n}} \neq p\left(y_{n} \mid w_{n}, \beta\right)
$$

on a set with positive probability. Note that $E l_{n}(b)$ is concave so that it is uniquely minimized at $b=\beta$.

## Proof of Theorem 3:

For convergence in probability, we check the conditions of Theorem 2.7 of Newey and McFadden (1994). The objective function $L_{n}(b)$ is concave. The stationarity and strong mixing assumptions imply ergodicity, see White (2001, theorem 3.34). This implies pointwise convergence, $L_{n}(b) \xrightarrow{p} E l_{n}(b)$ for all $b$. Lemma 1 proves that $E l_{n}(b)$ is uniquely maximized at $\beta$. Therefore, all conditions of Theorem 2.7 of Newey and McFadden (1994) are satisfied and consistency follows. For almost sure convergence, note that it is easily seen from Lemma 1 and Lemma 2 that all the conditions of Theorem A1 of Wooldridge (1994) are satisfied, except for the condition of uniform convergence in probability of $L_{N}(b)$. Note that Wooldridge's Theorem A1 can be extended to include a strong convergence result if instead of uniform convergence in probability of $L_{N}(b)$, uniform almost sure convergence $L_{N}(b)$ is assumed. To show this uniform convergence, we use the generic uniform law of large numbers of Andrews (1987). To show strong uniform law of large numbers, this theorem requires compactness of the parameter space, and in addition it needs to be verified that the summands $q_{n}\left(w_{n}, b\right)$ are such that $q_{n}\left(w_{n}, b\right), q_{n}^{*}\left(w_{n}, b\right)=\sup \left\{q_{n}\left(w_{n}, \tilde{b}\right): \tilde{b} \in B,|\tilde{b}-b|<\rho\right\}$ and $q_{n *}\left(w_{n}, b\right)=\inf \left\{q_{n}\left(w_{n}, \tilde{b}\right): b \in B,|\tilde{b}-b|<\rho\right\}$ are well-defined and satisfy a strong law of large numbers, and that for all $b \in B$,

$$
\lim _{\rho \rightarrow 0} \sup _{n \in \mathbb{Z}}\left|N^{-1} \sum_{n=1}^{N} E q_{n}^{*}\left(w_{n}, b\right)-E q_{* n}\left(w_{n}, b\right)\right|=0
$$

The latter condition follows from stationarity of $\left(y_{n}, x_{n}\right)$, continuity, and the envelope condition of Assumption A. In addition, $q_{n}\left(w_{n}, b\right), q_{n}^{*}\left(w_{n}, b\right)$ and $q_{n *}\left(w_{n}, b\right)$ are well-defined and strong mixing random variables, so that we can apply the strong law of large numbers of Theorem 4 of de Jong (1995), from which it follows that if $\alpha(m)+\nu(m) \leq C m^{-\eta}$ for some positive constants $C$ and $\eta$, these variables will satisfy a strong law of large numbers. This is because under the condition that $E\left|l_{n}(b)\right|^{1+\delta}<\infty$, the summands will be an $L_{1+\delta / 2}$-mixingale.

Lemma 3 Under the conditions of Theorem 4,

$$
\left.(N-p)^{1 / 2}\left(\partial L_{N}(b) / \partial b\right)\right|_{b=\beta} \xrightarrow{d} N(0, I) .
$$

## Proof:

Note that by assumption, $E\left(\left.\left(\partial L_{n}(b) / \partial b\right)\right|_{b=\beta} \mid w_{n}\right)=0$ so that $\left.E\left(\partial L_{n}(b) / \partial b\right)\right|_{b=\beta}=0$. Moreover, $\left.\left(\partial L_{n}(b) / \partial b\right)\right|_{b=\beta}$ is a martingale difference sequence that is strong mixing and strictly stationary. In particular, the version of Bierens (2004, Theorem 7.11) of a central limit theorem of McLeish (1974) yields asymptotic normality. Applying the information matrix equality yields the result.

## Proof of Theorem 4:

We prove Theorem 4 by checking the conditions of Newey and McFadden (1994, theorem 3.1). Consistency was shown in Theorem 3. Condition (i) was assumed. Condition (ii), twice differentiability of the log likelihood, follows from the probit specification. Condition (iii) was shown in Lemma 3. Note that stationarity and strong mixing imply ergodicity, see White (2001, theorem 3.34). Condition (iv) then follows from the probit specification and reasoning similar to Newey and McFadden, page 2147, example 1.2. Nonsingularity follows from the probit specification and $E w_{n} w_{n}^{\prime}$ being positive definite so that condition (iv) is satisfied.

For the proof of Theorem 5, we need the following lemmas.

Lemma 4 For all $a \in \mathbb{R}$, if $0 \leq z_{n} \leq 1$ and $\left(z_{n}, x_{n}\right)$ is strictly stationary and strong mixing, then

$$
\sup _{b \in B}\left|N^{-1} \sum_{n=1}^{N}\left(z_{n} I\left(b^{\prime} x_{n} \leq a\right)-E z_{n} I\left(b^{\prime} x_{n} \leq a\right)\right)\right| \xrightarrow{p} 0 .
$$

In addition, if $\alpha(m) \leq C m^{-\eta}$ for positive constants $C$ and $\eta$, the convergence is almost surely.

## Proof of Lemma 4:

We will apply the generic uniform law of large numbers of the Theorem of Andrews (1987). It requires compactness of the parameter space $B$ (which is assumed), and in addition it needs to be verified that the summands $q_{n}\left(w_{n}, b\right)$ are such that $q_{n}\left(w_{n}, b\right), q_{n}^{*}\left(w_{n}, b\right)=\sup \left\{q_{n}\left(w_{n}, \tilde{b}\right)\right.$ : $\tilde{b} \in B,|\tilde{b}-b|<\rho\}$ and $q_{n *}\left(w_{n}, b\right)=\inf \left\{q_{n}\left(w_{n}, \tilde{b}\right): b \in B,|\tilde{b}-b|<\rho\right\}$ are well-defined and satisfy a (respectively weak or strong) law of large numbers, and for all $b \in B$,

$$
\lim _{\rho \rightarrow 0} \sup _{n \in \mathbb{Z}}\left|N^{-1} \sum_{n=1}^{N} E q_{n}^{*}\left(w_{n}, b\right)-E q_{* n}\left(w_{n}, b\right)\right|=0
$$

To show the last result, note that $\left(z_{n}, x_{n}\right)$ is strictly stationary under the conditions of the theorem, and therefore

$$
\begin{aligned}
& \lim _{\rho \rightarrow 0} \sup _{n \in \mathbb{Z}}\left|N^{-1} \sum_{n=1}^{N} E q_{n}^{*}\left(w_{n}, b\right)-E q_{* n}\left(w_{n}, b\right)\right| \\
& =\lim _{\rho \rightarrow 0} \sup _{n \in \mathbb{Z}}\left|E z_{n} I\left(\sup _{\tilde{b}:|b-\tilde{b}|<\rho} b^{\prime} x_{n}<a\right)-E z_{n} I\left(\inf _{\tilde{b}:|b-\tilde{b}|<\rho} b^{\prime} x_{n}<a\right)\right| \\
& \leq \limsup _{K \rightarrow \infty} \lim _{\rho \rightarrow 0} \sup _{n \in \mathbb{Z}} \mid E z_{n}\left(I\left(b^{\prime} x_{n}<a+\rho\left|x_{n}\right|\right)-I\left(b^{\prime} x_{n}<a-\rho\left|x_{n}\right|\right)\right) I\left(\left|x_{n}\right| \leq K\right) \\
& \left.\quad \quad+\limsup _{K \rightarrow \infty} \lim _{\rho \rightarrow 0} \sup _{n \in \mathbb{Z}} \mid E z_{n}\left(I\left(b^{\prime} x_{n}<a+\rho\left|x_{n}\right|\right)-I\left(b^{\prime} x_{n}<a-\rho\left|x_{n}\right|\right)\right) I\left(\left|x_{n}\right|>K\right)\right) \mid \\
& \leq \limsup _{K \rightarrow \infty} \lim _{\rho \rightarrow 0}\left|P\left(b^{\prime} x_{n}<a+\rho K\right)-P\left(b^{\prime} x_{n}<a-\rho K\right)\right|+\limsup _{K \rightarrow \infty} P\left(\left|x_{n}\right|>K\right)=0,
\end{aligned}
$$

because $x_{1 n}$ has a continuous distribution. Furthermore, note that $q_{n}\left(z_{n}, b\right)$,

$$
q_{n}^{*}\left(w_{n}, b\right)=z_{n} I\left(\sup _{\tilde{b}:|b-\tilde{b}|<\rho} b^{\prime} x_{n}<a\right)
$$

and

$$
q_{* n}\left(w_{n}, b\right)=z_{n} I\left(\inf _{\tilde{b}:|b-\tilde{b}|<\rho} b^{\prime} x_{n}<a\right)
$$

are well-defined and strong mixing random variables, implying that weak law of large numbers for mixingales of Andrews (1988) applies; or alternatively we can apply the strong law of large numbers of Theorem 4 of de Jong (1995), from which it follows that if $\alpha(m)+\nu(m) \leq$ $C m^{-\eta}$ for some positive constants $C$ and $\eta$, these variables will satisfy a strong law of large numbers (note that because of boundedness of the summands, the summands are $L_{2^{-}}$ mixingales).

Lemma 5 Under Assumptions 1,2, 3, 4 and 5,

$$
\sup _{b \in B}\left|S_{N}\left(b, \sigma_{N}\right)-E S_{N}\left(b, \sigma_{N}\right)\right| \xrightarrow{p} 0 .
$$

In addition, if $\alpha(m) \leq C m^{-\eta}$ for positive constants $C$ and $\eta$, the convergence is almost surely.

## Proof of Lemma 5:

First note that Horowitz' proof of his Lemma 4 (i.e. $\sup _{b \in B}\left|S_{N}\left(b, \sigma_{N}\right)-S_{N}^{*}(b)\right| \xrightarrow{a s} 0$ ) goes through as it stands, except for the proof of uniform convergence of the term in his Equation (A4), which uses a uniform law of large numbers for i.i.d. random variables. To show that

$$
\sup _{b \in B}\left|N^{-1} \sum_{n=1}^{N}\left(I\left(\left|\sum_{j=1}^{p} r_{j} y_{n-j}+c^{\prime} x_{n}\right|<\alpha\right)-E I\left(\left|\sum_{j=1}^{p} r_{j} y_{n-j}+c^{\prime} x_{n}\right|<\alpha\right)\right)\right|
$$

satisfies a strong or weak law of large numbers, we can use Lemma 4 . To do so, note that

$$
\begin{aligned}
& N^{-1} \sum_{n=1}^{N}\left(I\left(\left|\sum_{j=1}^{p} r_{j} y_{n-j}+c^{\prime} x_{n}\right|<\alpha\right)\right. \\
& =\sum_{j_{1}=0}^{1} \ldots \sum_{j_{p}=0}^{1} N^{-1} \sum_{n=1}^{N} I\left(y_{n-1}=j_{1}\right) \ldots I\left(y_{n-p}=j_{p}\right) I\left(\left|\sum_{i=1}^{p} r_{i} j_{i}+c^{\prime} x_{n}\right|<\alpha\right)
\end{aligned}
$$

and note that $I\left(y_{n-1}=j_{1}\right) \ldots I\left(y_{n-p}=j_{p}\right)$ is strong mixing, because it is the product of strong mixing random variables. It now only remains to be proven that

$$
\sup _{b \in B}\left|S_{N}^{*}(b)-S(b)\right| \xrightarrow{p} 0 \quad \text { or } \quad \xrightarrow{a s} 0,
$$

which Horowitz shows by referring to Manski (1985). This can be shown by noting that

$$
\begin{aligned}
& S_{N}^{*}(b)=N^{-1} \sum_{n=1}^{N}\left(2 \cdot I\left(y_{n}=1\right)-1\right) I\left(b^{\prime} x_{n} \geq 0\right) \\
& =2 N^{-1} \sum_{n=1}^{N} y_{n} I\left(b^{\prime} x_{n} \geq 0\right)-N^{-1} \sum_{n=1}^{N} I\left(b^{\prime} x_{n} \geq 0\right)
\end{aligned}
$$

and by Lemma 4, both terms satisfy a (weak or strong) uniform law of large numbers.

Lemma 6 Under Assumptions 1,2, 3 and 4, $S(b) \leq S(\beta)$ with equality holding only if $b=\beta$.

## Proof of Lemma 6:

This result follows by noting that all conditions from Lemma 3 of Manski (1985) are satisfied.

## Proof of Theorem 5 :

The proof of the theorem now follows from Theorem A1 of Wooldridge (1994) and the results of Lemma 5 and Lemma 6.

Let $z_{n}=\sum_{j=1}^{p} \rho_{j} y_{n-j}+\gamma^{\prime} x_{n}$. The following lemma shows that Horowitz' Lemma 5 holds as it stands in our setting:

Lemma 7 Under Assumptions 1' and Assumptions 2-14,

$$
\begin{aligned}
& \lim _{N \rightarrow \infty} E\left[\sigma_{N}^{-h} T_{N}\left(\beta, \sigma_{N}\right)\right]=A \\
& \lim _{N \rightarrow \infty} \operatorname{Var}\left[\left(N \sigma_{N}\right)^{1 / 2} T_{N}\left(\beta, \sigma_{N}\right)\right]=D .
\end{aligned}
$$

## Proof of Lemma 7:

The only adjustment to Horowitz' Lemma 5 that needs to be made is to show that the covariance terms in $\operatorname{Var}\left[\left(N \sigma_{N}\right)^{1 / 2} T_{N}\left(\beta, \sigma_{N}\right)\right]$ are asymptotically negligible. To prove this, we show that for all vectors $\xi$ such that $|\xi|=1$,

$$
\lim _{N \rightarrow \infty} \sigma_{N} \sum_{m=1}^{\infty}\left|\operatorname{cov}\left(\xi^{\prime}\left(\tilde{x}_{n} / \sigma_{N}\right) K^{\prime}\left(z_{n} / \sigma_{N}\right), \xi^{\prime}\left(\tilde{x}_{n-m} / \sigma_{N}\right) K^{\prime}\left(z_{n-m} / \sigma_{N}\right)\right)\right|=0
$$

By the covariance inequality for mixingales, for the same $s$ as in Assumption 6, (see Davidson (1994, p. 212, Corollary 14.3)),

$$
\begin{aligned}
& \sigma_{N} \operatorname{cov}\left(\xi^{\prime}\left(\tilde{x}_{n} / \sigma_{N}\right) K^{\prime}\left(z_{n} / \sigma_{N}\right), \xi^{\prime}\left(\tilde{x}_{n-m} / \sigma_{N}\right) K^{\prime}\left(z_{n-m} / \sigma_{N}\right)\right) \\
& \leq \sigma_{N} C \alpha(m)^{1-2 / s}\left(E\left|\xi^{\prime}\left(\tilde{x}_{n} / \sigma_{N}\right) K^{\prime}\left(z_{n} / \sigma_{N}\right)\right|^{s}\right)^{1 / s}\left(E\left|\xi^{\prime}\left(\tilde{x}_{n-m} / \sigma_{N}\right) K^{\prime}\left(z_{n-m} / \sigma_{N}\right)\right|^{s}\right)^{1 / s} \\
& =\sigma_{N}^{-1} C \alpha(m)^{1-2 / s}\left(\int\left|\xi^{\prime} \tilde{x}\right|^{s}\left|K^{\prime}\left(z / \sigma_{N}\right)\right|^{s} p(z \mid \tilde{x}) d z d P(\tilde{x})\right)^{2 / s} \\
& =C \alpha(m)^{1-2 / s} \sigma_{N}^{2 / s-1}\left(\int\left|\xi^{\prime} \tilde{x}\right|^{s}\left|K^{\prime}(\zeta)\right|^{s} p\left(\sigma_{N} \zeta \mid \tilde{x}\right) d \zeta d P(\tilde{x})\right)^{2 / s}
\end{aligned}
$$

by substituting $\zeta=z / \sigma_{N}$. The last term is smaller than $C^{\prime} \sigma_{N}^{2 / s-1} \alpha(m)^{1-2 / s}$ for some constant $C^{\prime}$. In view of the fact that summing the latter expression over $m$ will give a term that diverges as $N \rightarrow \infty$, we also need to use a second bound. To obtain this second bound, note that by Horowitz' arguments, under the conditions of the theorem,

$$
\sigma_{N} E \xi^{\prime}\left(\tilde{x}_{n-m} / \sigma_{N}\right) K^{\prime}\left(z_{n-m} / \sigma_{N}\right)=O\left(\sigma_{N}\right),
$$

implying that

$$
\begin{aligned}
& \quad \sigma_{N} \operatorname{cov}\left(\xi^{\prime}\left(\tilde{x}_{n} / \sigma_{N}\right) K^{\prime}\left(z_{n} / \sigma_{N}\right), \xi^{\prime}\left(\tilde{x}_{n-m} / \sigma_{N}\right) K^{\prime}\left(z_{n-m} / \sigma_{N}\right)\right) \\
& =O\left(\sigma_{N}\right)+\sigma_{N} E\left(\sigma_{N} \xi^{\prime}\left(\tilde{x}_{n} / \sigma_{N}\right) K^{\prime}\left(z_{n} / \sigma_{N}\right) \xi^{\prime}\left(\tilde{x}_{n-m} / \sigma_{N}\right) K^{\prime}\left(z_{n-m} / \sigma_{N}\right)\right) \\
& =O\left(\sigma_{N}\right)+\sigma_{N}^{-1} \int \xi^{\prime} \tilde{x}_{n} K^{\prime}\left(z_{n} / \sigma_{N}\right) \xi^{\prime} \tilde{x}_{n-m} K^{\prime}\left(z_{n-m} / \sigma_{N}\right) d P\left(x_{n}, x_{n-m}, z_{n}, z_{n-m}\right) \\
& =O\left(\sigma_{N}\right)+\sigma_{N}^{-1} \int \xi^{\prime} \tilde{x}_{n} K^{\prime}\left(z_{n} / \sigma_{N}\right) \xi^{\prime} \tilde{x}_{n-m} K^{\prime}\left(z_{n-m} / \sigma_{N}\right) d p\left(z_{n}, z_{n-m} \mid x_{n}, x_{n-m}\right) d z_{n} d z_{n-m} d P\left(x_{n}, x_{n-m}\right) \\
& =O\left(\sigma_{N}\right)+\sigma_{N} \iint K^{\prime}\left(\zeta_{n}\right) K^{\prime}\left(\zeta_{n-m}\right) p\left(\sigma_{N} \zeta_{n}, \sigma_{N} \zeta_{n-m} \mid x_{n}, x_{n-m}\right) d \zeta_{n} d \zeta_{n-m} \xi^{\prime} \tilde{x}_{n-m} \xi^{\prime} \tilde{x}_{n} d P\left(x_{n}, x_{n-m}\right)
\end{aligned}
$$

$$
=O\left(\sigma_{N}\right)
$$

under the assumptions of the theorem. Therefore for any $\kappa \in(0,1)$,

$$
\begin{aligned}
& \sum_{m=1}^{\infty}\left|\operatorname{cov}\left(\xi^{\prime}\left(\tilde{x}_{n} / \sigma_{N}\right) K^{\prime}\left(z_{n} / \sigma_{N}\right), \xi^{\prime}\left(\tilde{x}_{n-m} / \sigma_{N}\right) K^{\prime}\left(z_{n-m} / \sigma_{N}\right)\right)\right| \\
& \leq C \sum_{m=1}^{\infty}\left(\sigma_{N}\right)^{\kappa}\left(\alpha(m)^{(1-2 / s)} \sigma_{N}^{2 / s-1}\right)^{1-\kappa}
\end{aligned}
$$

and by choosing $\kappa=(s-2) /(2 s-2)+\eta$ and $\eta>0$ small enough, the last term can be bounded by

$$
C\left(\sum_{m=1}^{\infty} \alpha(m)^{(s-2) /(2 s-2)-\eta(s-2) / s}\right) \sigma_{N}^{(2 s-2) \eta / s}=O\left(\sigma_{N}^{(2 s-2) \eta / s}\right)=o(1)
$$

where the finiteness of the summation follows from the assumptions.

Horowitz' Lemma 6 now holds as follows:
Lemma 8 Under Assumptions 1' and Assumptions 2-14, (a) If $N \sigma_{N}^{2 h+1} \rightarrow \infty$ as $N \rightarrow \infty$, $\sigma_{N}^{-h} T_{N}\left(\beta, \sigma_{N}\right) \xrightarrow{p}$ A. (b) If $N \sigma_{N}^{2 h+1}$ has a finite limit $\lambda$ as $N \rightarrow \infty,\left(N \sigma_{N}\right)^{1 / 2} T_{N}\left(\beta, \sigma_{N}\right) \xrightarrow{d}$ $N\left(\lambda^{1 / 2} A, D\right)$.

## Proof of Lemma 8:

The modification of Horowitz (1992) that is needed is to show that for all vectors $\xi$ such that $|\xi|=1$,

$$
\left(\sigma_{N} / N\right)^{1 / 2} \xi^{\prime} \sum_{n=1}^{N}\left(t_{N n}-E t_{N n}\right) \xrightarrow{d} N\left(0, \xi^{\prime} D \xi\right),
$$

where

$$
t_{N n}=\left(2 y_{n}-1\right)\left(\tilde{x}_{n} / \sigma_{N}\right) K^{\prime}\left(z_{n} / \sigma_{N}\right)
$$

Since $t_{N n}$ is strong mixing, Theorem 2 of de Jong (1997) for strong mixing arrays can now be applied to show this result under the conditions of the lemma. Note that the condition $\alpha(m) \leq C m^{-s /(s-2)-\eta}$ from that theorem follows from the assumptions of the lemma.

For reproving Horowitz' Lemma 7 for the case of strong mixing data, we need the following lemmas:

Lemma 9 (Azuma(1967)) If $\eta_{n}$ is a martingale difference sequence with respect to $\mathcal{F}_{n}$ and $\left|\eta_{n}\right| \leq C_{N}$, then

$$
P\left(\left|N^{-1} \sum_{n=1}^{N} \eta_{n}\right|>\delta\right) \leq 2 \exp \left(-N \delta^{2} / C_{N}^{2}\right)
$$

## Proof of Lemma 9:

See Azuma (1967).

An $m_{N}$-fold application of the above lemma now gives the following result:
Lemma 10 If $\mathcal{F}_{n}$ is a sequence of sigma-fields such that $\eta_{n}-E\left(\eta_{n} \mid \mathcal{F}_{n-1}\right)$ is a martingale difference sequence with respect to $\mathcal{F}_{n}$ and $\left|\eta_{n}\right| \leq C_{N}$, then for any integer-valued sequence $m_{N}$ such that $m_{N} \geq 1$,

$$
P\left(\left|N^{-1} \sum_{n=1}^{N}\left(\eta_{n}-E\left(\eta_{n} \mid \mathcal{F}_{n-m_{N}}\right)\right)\right|>\delta\right) \leq 2 m_{N} \exp \left(-\delta^{2} /\left(m_{N}^{2} C_{N}^{2}\right)\right)
$$

## Proof of Lemma 10:

Obviously

$$
N^{-1} \sum_{n=1}^{N}\left(\eta_{n}-E\left(\eta_{n} \mid \mathcal{F}_{n-m_{N}}\right)\right)=\sum_{j=0}^{m_{N}-1} N^{-1} \sum_{n=1}^{N}\left(E\left(\eta_{n} \mid \mathcal{F}_{n-j}\right)-E\left(\eta_{n} \mid \mathcal{F}_{n-j-1}\right)\right),
$$

and therefore

$$
\begin{aligned}
& P\left(\left|N^{-1} \sum_{n=1}^{N}\left(\eta_{n}-E\left(\eta_{n} \mid \mathcal{F}_{n-m_{N}}\right)\right)\right|>\delta\right) \\
& \leq \sum_{j=0}^{m_{N}-1} P\left(\left|N^{-1} \sum_{n=1}^{N}\left(E\left(\eta_{n} \mid \mathcal{F}_{n-j}\right)-E\left(\eta_{n} \mid \mathcal{F}_{n-j-1}\right)\right)\right|>\delta / m_{N}\right) \\
& \leq 2 m_{N} \exp \left(-\delta^{2} /\left(m_{N}^{2} C_{N}^{2}\right)\right)
\end{aligned}
$$

by $m_{N}$ applications of Lemma 9 .

Analogously to Horowitz (1992), define

$$
g_{N n}(\theta)=\left(2 \cdot I\left(y_{n}=1\right)-1\right) \tilde{x}_{n} K^{\prime}\left(z_{n} / \sigma_{N}+\theta^{\prime} \tilde{x}_{n}\right)
$$

The following result is now the analogue ${ }^{2}$ of Horowitz' Lemma 7.
Lemma 11 If $\left(y_{n}, x_{n}\right)$ is strong mixing with strong mixing sequence $\alpha(m)$, and there exists a sequence $m_{N} \geq 1$ such that

$$
\sigma_{N}^{-3(p+q-1)} \sigma_{N}^{-2} N^{1 / s} \alpha\left(m_{N}\right)+\left(\log \left(N m_{N}\right)\right)\left(N^{1-2 / s} \sigma_{N}^{4} m_{N}^{-2}\right)^{-1} \rightarrow 0,
$$

then

$$
\sup _{\theta \in \Theta_{N}}\left|\left(N \sigma_{N}^{2}\right)^{-1} \sum_{n=1}^{N}\left(g_{N n}(\theta)-E g_{N n}(\theta)\right)\right| \xrightarrow{p} 0 .
$$

Note that the second part of Horowitz' Lemma 7 will hold without modification. Also note that the case of i.i.d. $\left(y_{n}, x_{n}\right)$ is a special case, because then $\alpha(m)=0$ for $m \geq 1$, and we could set $m_{N}=1$ for that case.

## Proof of Lemma 11:

Consider

$$
g_{N n}^{C_{N}}(\theta)=\left(2 \cdot I\left(y_{n}=1\right)-1\right) \tilde{x}_{n} K^{\prime}\left(z_{n} / \sigma_{N}+\theta^{\prime} \tilde{x}_{n}\right) I\left(\left|\tilde{x}_{n}\right| \leq C_{N}\right)
$$

and note that obviously,

$$
g_{N n}(\theta)-E g_{N n}(\theta)=\left(g_{N n}^{C_{N}}(\theta)-E g_{N n}^{C_{N}}(\theta)\right)
$$

[^2]\[

$$
\begin{equation*}
+\left(g_{N n}(\theta)-g_{N n}^{C_{N}}(\theta)-E g_{N n}(\theta)+E g_{N n}^{C_{N}}(\theta)\right) \tag{41}
\end{equation*}
$$

\]

Now define $C_{N}=\eta^{-1 / s} N^{1 / s}\left(E\left|\tilde{x}_{n}\right|^{s}\right)^{1 / s}$ for any $\eta>0$. Then because $C_{N} \rightarrow \infty$ as $N \rightarrow \infty$, following the reasoning as in the proof of (A16) of Horowitz (1992, page 525-526), it follows that

$$
\begin{equation*}
\sup _{\theta \in \Theta}\left|E g_{N n}(\theta)-E g_{N n}^{C_{N}}(\theta)\right| \rightarrow 0 \tag{42}
\end{equation*}
$$

In addition,

$$
\begin{equation*}
P\left(\sup _{\theta \in \Theta}\left|\sum_{n=1}^{N}\left(g_{N n}(\theta)-g_{N n}^{C_{N}}(\theta)\right)\right|=0\right) \leq P\left(\exists n:\left|\tilde{x}_{n}\right|>C_{N}\right) \leq N E\left|\tilde{x}_{n}\right|^{s} C_{N}^{-s} \leq \eta \tag{43}
\end{equation*}
$$

and we can choose $\eta$ arbitrarily small. For the case $s=\infty$, it is trivial that these two terms disappear asymptotically for some constant $C_{N}$ not depending on $N$. To deal with the first part of Equation (41), note that

$$
\begin{equation*}
g_{N n}(\theta)-E g_{N n}(\theta)=\left(g_{N n}(\theta)-E\left(g_{N n}(\theta) \mid \mathcal{F}_{n-m_{N}}\right)\right)+\left(E\left(g_{N n}(\theta) \mid \mathcal{F}_{n-m_{N}}\right)-E g_{N n}(\theta)\right) \tag{44}
\end{equation*}
$$

To deal with the first part of the right-hand side of Equation (44), we can copy the argument on page 525 of Horowitz (1992), except that now, by Lemma 10,

$$
\begin{aligned}
& \sum_{i=1}^{\Gamma_{N}} P\left(\left(N \sigma_{N}^{2}\right)^{-1}\left|\sum_{n=1}^{N}\left(g_{N n}\left(\theta_{N i}\right)-E g_{N n}\left(\theta_{N i}\right)\right)\right|>\varepsilon / 2\right) \\
& \leq 2 \Gamma_{N} m_{N} \exp \left(-\varepsilon^{2} 4^{-1} N \sigma_{N}^{4} C_{N}^{-2} m_{N}^{-2}\right)
\end{aligned}
$$

where $\Gamma_{N}$ is as defined in Horowitz (1992). Since $\Gamma_{N}=O\left(\sigma_{N}^{-3(p+q-1)}\right)$, this term will converge to zero if

$$
\begin{equation*}
\left(\log \left(N m_{N}\right)\right)\left(N \sigma_{N}^{4} C_{N}^{-2} m_{N}^{-2}\right)^{-1} \rightarrow 0 \tag{45}
\end{equation*}
$$

which is assumed. For dealing with the second part of the right-hand side of Equation (44), note since $g_{N n}(\theta)$ is strong mixing, it is also an $L_{1}$-mixingale (see for example Davidson (1994, p. 249, Example 16.3), implying that

$$
E\left|E\left(g_{N n}^{C_{N}}(\theta) \mid \mathcal{F}_{n-m_{N}}\right)-E g_{N n}^{C_{N}}(\theta)\right| \leq 6 C_{N} \alpha\left(m_{N}\right)
$$

Using Horowitz' reasoning of page 525 , it now suffices to show that for all $\varepsilon>0$,

$$
\sum_{i=1}^{\Gamma_{N}} P\left(\left(N \sigma_{N}^{2}\right)^{-1}\left|\sum_{n=1}^{N} E\left(g_{N n}^{C_{N}}\left(\theta_{N i}\right) \mid \mathcal{F}_{n-m_{N}}\right)-E g_{N n}^{C_{N}}\left(\theta_{N i}\right)\right|>\varepsilon\right) \rightarrow 0
$$

By the Markov inequality,

$$
\begin{aligned}
& \sum_{i=1}^{\Gamma_{N}} P\left(\left(N \sigma_{N}^{2}\right)^{-1}\left|\sum_{n=1}^{N} E\left(g_{N n}^{C_{N}}\left(\theta_{N i}\right) \mid \mathcal{F}_{n-m_{N}}\right)-E g_{N n}^{C_{N}}\left(\theta_{N i}\right)\right|>\varepsilon\right) \\
& \leq \sum_{i=1}^{\Gamma_{N}} \varepsilon^{-1} \sigma_{N}^{-2} N^{-1} \sum_{n=1}^{N} E\left|E\left(g_{N n}^{C_{N}}(\theta) \mid \mathcal{F}_{n-m_{N}}\right)-E g_{N n}^{C_{N}}(\theta)\right| \\
& =O\left(\sigma_{N}^{-3(p+q-1)} \sigma_{N}^{-2} C_{N} \alpha\left(m_{N}\right)\right)=o(1)
\end{aligned}
$$

by assumption.

Lemma 12 Under Assumptions 1' and Assumptions 2-14, $\left(\tilde{b}_{N}-\tilde{\beta}\right) / \sigma_{N} \xrightarrow{p} 0$.

## Proof of Lemma 12:

This follows from Lemma 11 and the reasoning ${ }^{3}$ of Horowitz' (1992) Lemma 8.

The following lemma corresponds ${ }^{4}$ to Horowitz' Lemma 9.
Lemma 13 Let $\left\{\beta_{N}\right\}=\left\{\beta_{N 1}, \tilde{\beta}_{N}\right\}$ be such that $\left(\beta_{N}-\beta\right) / \sigma_{N} \xrightarrow{p} 0$ as $N \rightarrow \infty$. Then under Assumptions 1' and Assumptions 2-14,

$$
Q_{N}\left(\beta_{N}, \sigma_{N}\right) \xrightarrow{p} Q .
$$

[^3]
## Proof of Lemma 13:

Remember that

$$
Q_{N}\left(\beta_{N}, \sigma_{N}\right)=\left[\sigma_{N}^{-2} N^{-1} \sum_{n=1}^{N}\left(2 y_{n}-1\right) \tilde{x}_{n} \tilde{x}_{n}^{\prime} K^{\prime \prime}\left(\left(\sum_{j=1}^{p} r_{j} y_{n-j}+c^{\prime} x_{n}\right) / \sigma_{N}\right)\right]_{b=\beta_{N}}
$$

Since $P\left(b_{1}=\beta_{1}\right) \rightarrow 1$ and by the assumption that $\left(\beta_{N}-\beta\right) / \sigma_{N} \xrightarrow{p} 0$ as $N \rightarrow \infty$, it suffices to show that for all $\eta>0$ and any vector $\xi$ such that $|\xi|=1$,

$$
\begin{align*}
\sup _{|\tilde{\theta}| \leq \eta} \mid & N^{-1} \sum_{n=1}^{N} r_{n N}(\tilde{\theta})-E r_{n N}(\tilde{\theta})\left|\equiv \sup _{|\tilde{\theta}| \leq \eta}\right| \sigma_{N}^{-2} N^{-1} \sum_{n=1}^{N}\left(2 y_{n}-1\right)\left(\xi^{\prime} \tilde{x}_{n}\right)^{2} K^{\prime \prime}\left(z_{n} / \sigma_{N}+\tilde{\theta}^{\prime} \tilde{x}_{n}\right) \\
& -E\left(2 y_{n}-1\right)\left(\xi^{\prime} \tilde{x}_{n}\right)^{2} K^{\prime \prime}\left(z_{n} / \sigma_{N}+\tilde{\theta}^{\prime} \tilde{x}_{n}\right) \mid \xrightarrow{p} 0 . \tag{46}
\end{align*}
$$

Note that Horowitz (1992) shows the continuity of $E r_{n N}(\tilde{\theta})$ in $\tilde{\theta}$ uniformly in $N$. To show the result of Equation (46), note that

$$
\begin{aligned}
& P\left(\sup _{|\tilde{\theta}| \leq \eta}\left|N^{-1} \sum_{n=1}^{N} r_{n N}(\tilde{\theta}) I\left(\left|r_{n N}(\tilde{\theta})\right|>C_{N}\right)\right|=0\right) \\
& \leq \sum_{n=1}^{N} P\left(\left(\xi^{\prime} \tilde{x}_{n}\right)^{2}>C_{N}\right) \leq N E\left|\xi^{\prime} \tilde{x}_{n}\right|^{s} C_{N}^{-s / 2}
\end{aligned}
$$

and the last term can be made smaller than $\varepsilon$ by choosing $C_{N}^{-s / 2}=N^{-1} \varepsilon\left(E\left|\xi^{\prime} \tilde{x}_{n}\right|^{s}\right)^{-1}$. In addition, it is easily verified that

$$
\sup _{|\tilde{\theta}| \leq \eta}\left|N^{-1} \sum_{n=1}^{N} E\left(r_{n N}(\tilde{\theta}) I\left(\left|r_{n N}(\tilde{\theta})\right|>C_{N}\right)\right)\right| \rightarrow 0 .
$$

Because of these two results, it suffices to show uniform convergence to zero in probability of

$$
\begin{aligned}
& R_{N}(\tilde{\theta})=N^{-1} \sum_{n=1}^{N}\left(r_{n N}(\tilde{\theta}) I\left(\left|r_{n N}(\tilde{\theta})\right| \leq C_{N}\right)-E r_{n N}(\tilde{\theta}) I\left(\left|r_{n N}(\tilde{\theta})\right| \leq C_{N}\right)\right. \\
& \left.\quad+C_{N} I\left(\left|r_{n N}(\tilde{\theta})\right|>C_{N}\right)-E C_{N} I\left(\left|r_{n N}(\tilde{\theta})\right|>C_{N}\right)\right) .
\end{aligned}
$$

Now note that since $\tilde{\theta} \in \mathbb{R}^{p+q-1}$, we can cover the parameter space $\{\tilde{\theta}:|\tilde{\theta}| \leq \eta\}$ with $O\left(\sigma_{N}^{-2(p+q-1) / \mu}\right)$ balls of size $\sigma_{N}^{2 / \mu}$ and with centers $\tilde{\theta}_{j}$. Now note that, by Assumption 14,

$$
\begin{aligned}
& \sup _{N \geq 1} E \sup _{\left|\tilde{\theta}-\tilde{\theta}^{\prime}\right|<\delta \sigma_{N}^{2 / \mu}}\left|R_{N}(\tilde{\theta})-R_{N}\left(\tilde{\theta}^{\prime}\right)\right| \\
& \leq \sup _{N \geq 1} E\left(\xi^{\prime} \tilde{x}_{n}\right)^{2} L \sup _{\left|\tilde{\theta}-\tilde{\theta}^{\prime}\right|<\delta \sigma_{N}^{2 / \mu}}\left|\tilde{\theta}-\tilde{\theta}^{\prime}\right|^{\mu} \sigma_{N}^{-2} \rightarrow 0 \quad \delta \rightarrow 0 .
\end{aligned}
$$

Using Lemma 11 and following the same reasoning as in the proof of that lemma, we can now argue

$$
\begin{aligned}
& \limsup _{n \rightarrow \infty} P\left(\sup _{|\tilde{\theta}| \leq \eta}\left|R_{N}(\tilde{\theta})-E R_{N}(\tilde{\theta})\right|>\varepsilon\right) \\
& \leq \limsup _{n \rightarrow \infty} P\left(\max _{j}\left|R_{N}\left(\tilde{\theta}_{j}\right)-E R_{N}\left(\tilde{\theta}_{j}\right)\right|>\varepsilon / 2\right) \\
& \leq \limsup _{n \rightarrow \infty} \sum_{j} P\left(\left|R_{N}\left(\tilde{\theta}_{j}\right)-E R_{N}\left(\tilde{\theta}_{j}\right)\right|>\varepsilon / 2\right) \\
& =O\left(\sigma_{N}^{-2(p+q-1) / \mu}\left[2 m_{N} \exp \left(-N \varepsilon^{2} /\left(4 \sigma_{N}^{4} C_{N}^{2} m_{N}^{2}\right)\right)+\varepsilon^{-1} C_{N} \alpha\left(m_{N}\right)\right]\right)
\end{aligned}
$$

and because $C_{N}=O\left(N^{2 / s}\right)$, the last term converges to 0 if

$$
\sigma_{N}^{-2(p+q-1) / \mu} N^{2 / s} \alpha\left(m_{N}\right)+\left(m_{N}^{-2} \sigma_{N}^{-4} N^{1-4 / s}\right)^{-1} \log \left(N m_{N}\right) \rightarrow 0
$$

which is assumed.

## Proof of Theorem 6:

This proof is identical to the proof of Horowitz' Theorem 2, where we need to use our Lemma 12 and Lemma 13 instead of Horowitz' Lemma 8 and Lemma 9.

## Proof of Theorem 7:

Part (a) now follows exactly ${ }^{5}$ as in Horowitz' proof of his Theorem 3, where our Lemma 12 and Lemma 13 replace Horowitz' Lemma 8 and Lemma 9. Part (c) follows from Lemma 13.

[^4]for all $\xi$ such that $|\xi|=1$. Under the conditions of our theorem, this result can be proven analogously to the proof of Lemma 13, using the same $C_{N}$ and ball size sequences. Note that $K^{\prime}($.$) is Lipschitz-continuous$ with $\mu=1$, since $K^{\prime \prime}($.$) is assumed to exist and to be uniformly bounded.$


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[^1]:    ${ }^{1}$ In addition, some corrections to Horowitz' proof of the validity of the smoothed maximum score procedure are provided.

[^2]:    ${ }^{2}$ Note that Horowitz' Lemma 7 only holds for bounded regressors, and that the truncation argument at the start of Lemma 8 appears to be in error. Horowitz does not explicitly consider the remainder statistic containing the summation elements for which $\left|\tilde{x}_{n}\right|$ exceeds $a$. Horowitz' Lemma 9 appears to have a similar problem in its proof. Therefore, Lemma 11 also serves to correct this aspect of Horowitz' proof. This is because the conditioning on the event $C_{\gamma}$ does not appear relevant; while Horowitz' $\tilde{x}$ stands for a random variable distributed identically to any $\tilde{x}_{n}$, the conditioning should be with respect to every $\tilde{x}_{n}, n=1, \ldots, N$, in order for this argument to work. However, unless $\tilde{x}_{n}$ is almost surely bounded, such a conditioning set $C_{\gamma}$ would depend on $N$, and will not have the desired property that $\lim \sup _{\gamma \rightarrow \infty} \lim \sup _{N \rightarrow \infty} P\left(C_{\gamma}\right)=0$.

[^3]:    ${ }^{3}$ See footnote 2.
    ${ }^{4}$ Note that Horowitz' conditioning on $X_{N}$ appears to be in error, and note that when Horowitz uses his Lemma 8 in the proof of his Theorem 2, a uniform law of large numbers appears to be needed rather than the result of his Lemma 8 .

[^4]:    ${ }^{5}$ To show part (b), one can use a uniform law of large numbers result of the type

    $$
    \sup _{|\tilde{\theta}| \leq \eta} \mid \sigma_{N}^{-1} N^{-1} \sum_{n=1}^{N}\left(\xi^{\prime} \tilde{x}_{n}\right)^{2} K^{\prime}\left(z_{n} / \sigma_{N}+\tilde{\theta}^{\prime} \tilde{x}_{n}\right)-E\left(\left(\xi^{\prime} \tilde{x}_{n}\right)^{2} K^{\prime}\left(z_{n} / \sigma_{N}+\tilde{\theta}^{\prime} \tilde{x}_{n}\right) \mid \xrightarrow{p} 0\right.
    $$

