The Asymptotic Size and Power of the Augmented Dickey-Fuller Test for a Unit Root

Efstathios Paparoditis∗ Dimitris N. Politis†

Abstract

It is shown that the limiting distribution of the augmented Dickey-Fuller (ADF) test under the null hypothesis of a unit root is valid under a very general set of assumptions that goes far beyond the linear AR(∞) process assumption typically imposed. In essence, all that is required is that the error process driving the random walk possesses a continuous spectral density that is strictly positive. Furthermore, under the same weak assumptions, the limiting distribution of the ADF test is derived under the alternative of stationarity, and a theoretical explanation is given for the well-known empirical fact that the test’s power is a decreasing function of the chosen autoregressive order $p$. The intuitive reason for the reduced power of the ADF test is that, as $p$ tends to infinity, the $p$ regressors become asymptotically collinear.

Key words: Autoregressive Representation, Hypothesis Testing, Random Walk.

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∗Dept. of Mathematics and Statistics, University of Cyprus, P.O.Box 20537, CY 1678 Nicosia, CYPRUS; stathisp@ucy.ac.cy
†Department of Mathematics, University of California—San Diego, La Jolla, CA 92093-0112, USA; dpolitis@ucsd.edu
1 Introduction

Testing for the presence of a unit root is a widely investigated problem in econometrics; cf. Hamilton (1994) or Patterson (2011) for extensive treatments of this topic. Given a stretch of time series observations \(X_1, X_2, \ldots, X_n\), one of the commonly used tests for the null hypothesis of a unit root, is the so-called augmented Dickey-Fuller (ADF) test. This test decides about the presence of a unit root in the data generating mechanism by using the ordinary least squares (OLS) estimator \(\hat{\rho}_n\) of \(\rho\), obtained by fitting the regression equation

\[
X_t = \rho X_{t-1} + \sum_{j=1}^{p} a_{j,p} \Delta X_{t-j} + e_{t,p},
\]

(1.1)

to the observed stretch of data. In the above notation, \(\Delta X_t = X_t - X_{t-1}\), while the order \(p\) is allowed to depend on \(n\), i.e., \(p\) is short-hand for \(p_n\), in a way that is related to the assumptions imposed on the underlying process. In particular, under the null hypothesis \(H_0: \rho = 1\), it is commonly assumed that \(X_t\) is obtained by integrating a linear, infinite order autoregressive process, \((\text{AR}(\infty))\), i.e., that

\[
X_t = X_{t-1} + U_t, \quad t = 1, 2, \ldots,
\]

(1.2)

where \(X_0 = 0\) and

\[
U_t = \sum_{j=1}^{\infty} a_j U_{t-j} + e_t.
\]

(1.3)

Here \(\{e_t\}\) is a sequence of independent, identically distributed (i.i.d.) random variables having mean zero and variance \(0 < \sigma_e^2 < \infty\). Stationarity and causality of \(\{U_t\}\) is ensured by assuming that \(\sum_{j=1}^{\infty} |j|^s |a_j| < \infty\) for some \(s \geq 1\) and \(\sum_{j=1}^{\infty} a_j z^j \neq 0\) for all \(|z| \leq 1\).

To test \(H_0\), Dickey and Fuller (1979) proposed the studentized statistic

\[
t_n = \frac{\hat{\rho}_n - 1}{\hat{\text{Std}}(\hat{\rho}_n)},
\]

(1.4)

where \(\hat{\text{Std}}(\hat{\rho}_n)\) denotes an estimator of the standard deviation of the OLS estimator \(\hat{\rho}_n\). The asymptotic distribution of \(t_n\) under \(H_0\) is non-standard and is well known in the
literature. Dickey and Fuller (1979) and Dickey and Fuller (1981) derived this distribution under the assumption that the order of the underlying autoregressive process is finite and known. Said and Dickey (1984) extended this result for the case where the innovation process \( \{U_t\} \) driving the random walk (1.2) is an invertible autoregressive moving-average (ARMA) process, i.e., an AR(\(\infty\)) process with exponentially decaying coefficients. Ng and Perron (1995) relaxed the assumptions needed on the rate at which the order \( p \) in (1.1) increases to infinity with \( n \). Chang and Park (2002) established the same limiting distribution of \( t_n \) by further relaxing the assumptions regarding the rate at which \( p \) increases to infinity, by allowing for a polynomial decrease of the coefficients \( a_j \) in the AR(\(\infty\)) representation (1.3) and by assuming a martingale difference structure instead of i.i.d. innovations \( e_t \), that is, by assuming that \( E(e_t|\mathcal{E}_{t-1}) = 0 \) and 
\[
 n^{-1} \sum_{t=1}^{n} E(e_t^2|\mathcal{E}_{t-1}) \to \sigma^2, \quad n \to \infty,
\]
where \( \mathcal{E}_l = \sigma(\{e_t : t \leq l\}) \) is the \( \sigma \)-algebra generated by the random variables \( \{e_t, e_{t-1}, \ldots\} \).

To derive the power behavior of the test under the alternative hypothesis \( H_1 : \rho < 1 \), the limiting distribution of \( t_n \) is required under the assumption that \( \{X_t\} \) is a stationary process. Investigating the power of the ADF-test for fixed (stationary) alternatives, has attracted less interest in the literature. Nabeya and Tanaka (1990) and Perron (1991) analyzed the limiting power of unit root tests for sequences of local alternatives. For a first order autoregression, Abadir (1993) gives closed forms for the distribution of certain statistics leading to the derivation of the limiting distribution of unit root tests under the null and the alternative. For the ADF unit root test, Lopez (1997) considered the asymptotic distribution of \( \hat{\rho}_n \) under the alternative that \( \{X_t\} \) is a causal and invertible ARMA process; nevertheless, apart from the restrictive ARMA set-up, there is a flaw in his derivations. To elaborate, Lopez (1997) replaces the regression equation (1.1) by a regression that contains only levels of the \( X_t \)’s; but this replacement is not justified when \( p \) is of larger order of magnitude than \( n^{1/8} \) as our Remark 2.4 explains.

Testing for the presence of a unit root remains an active area of econometric research. As regards the ADF test, Elliot et al. (1996) proposed nearly efficient tests
while Harvey et al. (1999) discussed issues caused by the uncertainty about the deterministic trend and the initial values. The effects of the initial value in unit root testing were definitively investigated in the influential paper by Müller and Elliot (2003). In addition, Müller (2005, 2007) considered strongly autocorrelated time series and studied size and power properties of unit root tests in a local-to-unity asymptotic framework. Notably, under unit root (or local-to-unity) assumptions, the initial value will have an effect on the asymptotic distribution of the test statistic only if it is of order of magnitude of at least $\sqrt{n}$. However, in the standard asymptotic framework (that we also adopt), the effect of a fixed initial value vanishes as $n \to \infty$.

A different and interesting direction of research in unit root testing is described by so-called ‘tuning parameters free’ tests. Such tests have been considered by Breitung (2002), Müller and Watson (2008) and Nielsen (2009). By means of an appropriate rescaling, these tests bypass the problem of estimating the long-run variance which appears as a nuisance parameter in unit root testing. In the context of the well-known Phillips and Perron test—see Phillips (1987) and Phillips and Perron (1988)—, this nuisance parameter is estimated via non-parametric spectral methods whereas, in the context of the ADF test, the same problem is addressed via augmenting the regression equation by $p$ lagged differences, and allowing for $p$ to increase with sample size. Finally, recall that the popular KPSS stationarity test of Kwiatkowski et al. (1992) can also be employed for unit root testing; see e.g. Shin and Schmidt (1992).

Despite these developments, our paper contributes towards a deeper understanding of the ADF test. Firstly, we show that the asymptotic distribution of the ADF test statistic under the null hypothesis of a unit root is valid under a most general set of assumptions regarding the innovation process $\{U_t\}$ driving the random walk (1.2). These assumptions go far beyond the AR($\infty$) linear process class (1.3). In particular, we prove validity of the limiting distribution of $t_n$ under the general condition that the stationary process $\{U_t\}$ possesses a Wold-type, AR-representation with respect to white noise errors $\varepsilon_t$; see eq. (2.6) in what follows. This much wider class of stationary
processes should not be confused with the linear AR(∞) class (1.3) driven by i.i.d. or by martingale difference innovations. In fact, the class of processes admitting a Wold-type, AR-representation with respect to white noise errors consists of all zero mean, second order stationary (linear or nonlinear) processes having a continuous and strictly positive spectral density; cf. Pourahmadi (2001) for details.

Secondly, under the same set of general assumptions on the underlying stationary process \( \{U_t\} \), we establish the limiting distribution of the ADF-test \( t_n \) under the alternative hypothesis in which \( \{X_t\} \) is stationary. It turns out that under the alternative, the estimator \( \hat{\rho}_n \) is only \( \sqrt{n}/p \)–consistent, and therefore, its convergence rate is considerably smaller compared to the \( n \)–consistency of the same estimator under the null, and to the \( \sqrt{n} \) convergence rate of other test statistics under the alternative as, e.g., the test of Philips and Perron (1988). We make the case that the underlying reason for the slow rate of convergence of \( \hat{\rho}_n \) under the alternative is that the regressors in equation (1.1) become asymptotically collinear as \( p \) diverges; this is a surprising and novel result in the literature of unit root tests. We analyze this phenomenon, and show that this collinearity problem is responsible for the reduced power of the ADF-test, and gives a theoretic basis for the empirically observed fact that the power of the ADF test is a decreasing function of the order \( p \)—see e.g. Figure 9.5 of Patterson (2011).

The remainder of the paper is organized as follows. Section 2 states the main assumptions imposed on the underlying process \( \{U_t\} \) and derives the asymptotic distribution of the ADF test \( t_n \) under the null hypothesis of a unit root. The asymptotic behavior of \( t_n \) under the alternative of stationarity is also derived in Section 2, and its consequences for the power properties of the ADF test are discussed. Section 3 show cases a real data example where the reduced power of the ADF test is manifested in an extreme fashion, and underscores the importance of properly choosing the order \( p \) in practice. Section 4 presents the results of a simulation experiment that gives empirical confirmation of some of the theoretical results derived in Section 2. All technical proofs are deferred to Section 5.
2 Asymptotic Properties of the ADF test

2.1 Assumptions

We first state the conditions we impose on the dependence structure of the underlying second order stationary process \( \{U_t\} \) that drives the random walk under the null. Assuming that \( \{U_t\} \) is purely non-deterministic, i.e., that it possesses as spectral density whose logarithm is integrable, then the Wold representation yields

\[
U_t = \sum_{j=1}^{\infty} \alpha_j \varepsilon_{t-j} + \varepsilon_t
\]

where \( \sum_{j=1}^{\infty} \alpha_j^2 < \infty \) and \( \varepsilon_t \) is a zero mean, uncorrelated process with \( 0 < \text{Var}(\varepsilon_t) = \sigma^2 < \infty \). We slightly restrict the above class of stationary processes to the one satisfying the following assumption.

Assumption 1

(i) The autocovariance function \( \gamma_U(h) = \text{Cov}(U_t, U_{t+h}), h \in \mathbb{Z} \), of \( U = \{U_t, t \in \mathbb{Z}\} \) satisfies \( \sum_{h \in \mathbb{Z}} |\gamma_U(h)| < \infty \), and the spectral density \( f_U \) of \( U \) is strictly positive, i.e., \( f_U(\lambda) > 0 \) for all \( \lambda \).

(ii) \( E(\varepsilon_t^4) < \infty \) and the process \( \{\varepsilon_t\} \) satisfies the following weak dependence condition:

\[
\sum_{n=1}^{\infty} n \left\| P_1(\varepsilon_n) \right\| < \infty,
\]

where \( P_t(Y) = E(Y|F_t) - E(Y|F_{t-1}) \) and \( F_s = (\ldots, \varepsilon_{s-1}, \varepsilon_s) \).

Here, and in the rest of the paper, the norm \( \| \cdot \| \) is taken to be the \( L_p \) norm with \( p = 4 \).

The Wold representation (2.5) is an MA(\( \infty \)) equation for \( U_t \) with respect to its innovations \( \varepsilon_t \). Assumption 1(i) allows for an alternative representation of \( U_t \) with respect to the same white noise given in (2.5). In fact, because of the summability of the autocovariance function, the process \( \{U_t\} \) has a continuous spectral density \( f_U(\lambda) = (2\pi)^{-1} \sum_{h \in \mathbb{Z}} \gamma_U(h) \cos(\lambda h) \). This together with the strict positivity of \( f_U(\lambda) \)
implies that \( \{U_t\} \) possesses a so-called \textit{Wold-type AR-representation}, that is, \( U_t \) can be expressed as
\[
U_t = \sum_{j=1}^{\infty} b_j U_{t-j} + \varepsilon_t,
\] (2.6)
where \( \varepsilon_t \) is the same white noise innovation process as the one appearing in the Wold representation (2.5). Furthermore, the coefficients \( b_j \) are absolutely summable, i.e.,
\[
\sum_{j=1}^{\infty} |b_j| < \infty \text{ and } b(z) = 1 - \sum_{k=1}^{\infty} b_j z^j \neq 0 \text{ for } |z| \leq 1; \text{ see Pourahmadi (2001), Lemma 6.4.}
\]

The Wold-type AR-representation (2.6) of \( U_t \) with respect to the white noise process \( \varepsilon_t \) should not be confused with the rather strong assumption of a linear AR(\( \infty \)) process with respect to i.i.d. innovations. For example, one important difference between the class of process obeying a Wold-type AR-representation and the class of linear AR(\( \infty \)) processes (1.3) is the linearity of the optimal predictor. In fact, for processes in the class (1.3) with i.i.d. or with martingale difference errors, the optimal \( k \)-step ahead predictor, is always the linear predictor. That is, for positive \( k \), the general \( L_2 \)-optimal predictor of \( U_{t+k} \) based on its past \( U_t \), i.e., the conditional expectation \( E(U_{t+k}|U_s, s \leq t) \), is for processes (1.3) with i.i.d. or with martingale difference innovations, identical to the best linear predictor \( \mathcal{P}_{\mathcal{M}_t}(U_{t+k}) \). Here \( \mathcal{P}_C(Y) \) denotes orthogonal projection of \( Y \) onto the set \( C \) and \( \mathcal{M}_t = \text{span}\{U_j : j \leq s\} \), i.e., the closed linear span generated by the random variables \( \{U_j : j \leq s\} \). This linearity property of the \( L_2 \)-optimal predictor is not shared by processes that only admit a Wold-type AR-representation with respect to white noise innovations.

It is apparent that the class of processes having a Wold-type AR-representation is very large, and includes basically all time series, linear or \textit{nonlinear}, as long as they possess a strictly positive and continuous spectral density. The difference between the Wold-type AR-representation and the linear AR(\( \infty \)) property (1.3) is further illustrated
by means of the following two examples.

**Example 1:** (Non causal linear processes) Consider the process $U_t = \phi U_{t-1} + e_t$, with $|\phi| > 1$, and $e_t$ a zero mean i.i.d. process with variance $\sigma^2_e$. Notice that $\{U_t\}$ is stationary but it does not belong to the linear AR($\infty$) class (1.3) since it is not causal (the root of $1 - \phi z = 0$ lies outside the unit disc). However, for $\varepsilon_t = \phi^{-2}(e_t - (\phi^2 - 1) \sum_{j \geq 1} \phi^{-j} e_{t+j})$, $U_t$ has the AR-representation $U_t = b U_{t-1} + \varepsilon_t$ with $b = 1/\phi$ and the (causal) Wold representation $U_t = \sum_{j=1}^{\infty} b^j \varepsilon_{t-j}$ with respect to the white noise process $\{\varepsilon_t\}$, where $E(\varepsilon_t) = 0$ and $Var(\varepsilon_t) = \phi^{-2} \sigma^2_e$.

**Example 2:** (Non invertible linear processes) Consider the process $U_t = \theta e_{t-1} + e_t$, with $|\theta| > 1$, and $e_t$ a zero mean i.i.d. process with variance $\sigma^2_e$. Notice that $\{U_t\}$ is stationary but it does not belong to the linear AR($\infty$) class (1.3) since it is not invertible (the root of $1 - \theta z = 0$ lies outside the unit disc). Now, for $\varepsilon_t = e_t + (\theta^{-1} - \theta) \sum_{j \geq 1} \theta^{-j+1} e_{t-j}$ we get $U_t = \varepsilon_t - \theta^{-1} \varepsilon_{t-1}$ and therefore, $U_t$ has the AR-representation $U_t = \sum_{j=1}^{\infty} b U_{t-j} + \varepsilon_t$, $b_j = -(1/\theta)^j$, with respect to the white noise process $\{\varepsilon_t\}$. Here $E(\varepsilon_t) = 0$ and $Var(\varepsilon_t) = \sigma^2_e(1 + \theta^2 - \theta^4)$.

Assumption 1(ii) is imposed in order to control the dependence structure of the innovation process $\{\varepsilon_t\}$ and consequently of $U_t$; cf. Wu and Min (2005). It is based on the concept of weak dependence introduced by Wu (2005) and allows together with (2.5) for a very broad class of possible processes. Wu and Min (2005) give many examples of processes belonging to this class, including many well-known *nonlinear* processes, e.g., ARCH/GARCH processes, threshold autoregressive processes, bilinear processes and random coefficient autoregressions.

Note that instead of the weak dependence assumption above, other measures could be also used to control the dependence structure of the innovations process $\{\varepsilon_t\}$ as well. For instance, the results presented in this paper can be derived also under the alternative assumption that the innovation process $\{\varepsilon_t\}$ in (2.5) is strong mixing with strong mixing coefficient $\alpha_m$ satisfying $\sum_{m=1}^{\infty} \alpha_m^{1/2} < \infty$. In any case, Assumption 1(ii)
extends considerably the class of stationary process allowed and goes far beyond the limited class of linear autoregressive processes with i.i.d. innovations or innovations that form martingale differences.

Based on Assumption 1 regarding the class of stationary process \( \{U_t\} \), Assumption 2 below specifies the generation mechanism of the underlying and observable process \( X = \{X_t, t \geq 0\} \).

**Assumption 2** The process \( X \) satisfies one (and only one) of the following two conditions:

(i) (Unit root case:) \( X_0 = 0 \) and \( X_t = X_{t-1} + U_t \) for \( t = 1, 2, \ldots \)

(ii) (Stationary case:) \( X_t = U_t \), for \( t = 0, 1, 2, \ldots \),

where \( \{U_t\} \) is the second order stationary process specified in Assumption 1.

The assumption \( X_0 = 0 \) simplifies notation and does not affect the asymptotic results of this paper for the unit root case. It can be replaced by other assumptions concerning the starting value \( X_0 \) provided this random variable remains bounded in probability. Assumption 2 simply states that \( X_t \) is either a stationary process satisfying Assumption 1 or it is obtained by integrating such a stationary process.

Notice that if Assumption 2 is true, then \( X_t \) obeys also the useful representation

\[
X_t = \rho X_{t-1} + \sum_{j=1}^{\infty} a_j \Delta X_{t-j} + \varepsilon_t, \quad (2.7)
\]

with \( \varepsilon_t \) the white noise process discussed in Assumption 1. To see this, notice that (2.7) is obviously true if Assumption 2(i) is satisfied with the choices \( \rho = 1 \) and \( \Delta X_{t-1} = X_t - X_{t-1} = U_t \). Furthermore, if Assumption 2(ii) is true then it is easily verified that \( X_t = (\sum_{j=1}^{\infty} b_j) X_{t-1} - \sum_{j=1}^{\infty} (\sum_{s=j+1}^{\infty} b_s) \Delta X_{t-j} + \varepsilon_t \), which implies that also in this case (2.7) is true with

\[
\rho = \sum_{j=1}^{\infty} b_j, \quad \text{and} \quad a_j = - \sum_{s=j+1}^{\infty} b_s, \quad j = 1, 2, \ldots
\]
Now, let
\[ \rho_{\text{min}} = \inf \left\{ \rho = \sum_{j=1}^{\infty} b_j : b_j, j = 1, 2, \ldots \text{ and } b(z) = 1 - \sum_{j=1}^{\infty} b_j z^j \neq 0 \text{ for } |z| \leq 1 \right\}. \]

The null and alternative hypothesis of interest can then be stated as
\begin{align*}
H_0 & : \rho = 1, & H_1 & : \rho \in (\rho_{\text{min}}, 1). \quad (2.8)
\end{align*}

Notice that \( H_0 \) is equivalent to Assumption 2(i) while \( H_1 \) to Assumption 2(ii). The range of values of \( \rho \) under the alternative \( H_1 \) is an interval since \( B = \{ b_j, j = 1, 2, \ldots : b(z) \neq 0 \text{ for } |z| \leq 1 \} \) is a convex set and the mapping \( g : B \to \mathbb{R} \) with \( g(b_1, b_2, \ldots) = \sum_{j=1}^{\infty} b_j \equiv \rho \) is continuous. Notice that \( \rho_{\text{min}} < -1 \) is also possible, for instance if \( U_t = \varepsilon_t - \theta \varepsilon_{t-1} \) with \( \theta \in (0.5, 1) \).

\textbf{Remark 2.1} It is common in the econometric literature to state Assumption 2 in the following different form:
\begin{equation}
X_t = aX_{t-1} + U_t, \quad (2.9)
\end{equation}

where \( \{ U_t \} \) is some zero mean, second order stationary process satisfying certain conditions. In this formulation, the case \( a = 1 \) is associated with the null hypothesis of unit root, while the case \( |a| < 1 \) with the alternative; see among others Ng and Perron (1995) and Chang and Park (2002). It is easily seen that the above formulation is a restatement of Assumption 2 in the sense that in both cases the same conditions are imposed on the underlying process \( \{ X_t \} \). If \( a = 1 \) this is obviously true while for \( |a| < 1 \), using the backshift operator \( L^s X_t = X_{t-s} \) we have that \( X_t = (1 - aL)^{-1} U_t = \sum_{j=0}^{\infty} a^j U_{t-j} \) and \( \{ X_t \} \) is stationary. However, if (2.9) is considered as a model for \( X_t \), then identifiability and interpretability problems occur for the parameter \( a \) unless, of course, \( a = 1 \). To see why, let \( X_t \) be a stationary series, and define the new stationary series \( V_t = X_t - bX_{t-1} \) where \( b \) is arbitrary; i.e., in the stationary case, eq. (2.9) holds true for \emph{any value} of the parameter \( a \) as long as it is not one. To make the parameter \( a \) identifiable in the stationary case, an additional condition must be imposed, e.g., that the series \( U_t \) is the innovation series of \( X_t \). Our Assumption 2 avoids these difficulties.
2.2 Limiting Distribution under the Null

The following theorem establishes the limiting distribution of the test statistic $t_n$ under $H_0$ in (2.8). It shows that this limiting distribution is identical to that obtained under the AR($\infty$) linearity or weak linearity assumption for $\{U_t\}$; cf. Dickey and Fuller (1981) and Chang and Park (2002).

**Theorem 2.1** Let Assumption 1 and Assumption 2(i) be satisfied and suppose that $p_n \to \infty$ as $n \to \infty$ such that $p_n/\sqrt{n} \to 0$. Then

$$t_n \Rightarrow \int_0^1 W(t)dW(t)/(\int_0^1 W^2(t)dt)^{1/2},$$

where $\{W(t), t \in [0, 1]\}$ is the standard Wiener process on $[0, 1]$.

By the above theorem, an asymptotic $\alpha$-level test of the null hypothesis of a unit root is given by rejecting $H_0$ whenever $t_n$ is smaller than $C_\alpha$, where $C_\alpha$ is the lower $\alpha$-percentage point of the distribution of $\int_0^1 W(t)dW(t)/\sqrt{\int_0^1 W^2(t)dt}$. Notice that since the class of stationary processes satisfying Assumption 1 is very rich and contains as special case many linear and nonlinear processes including the commonly used linear AR($\infty$) process driven by i.i.d. innovations or by martingale differences, Theorem 2.1 generalizes considerably previous results regarding the limiting distribution of the ADF-test under the null hypothesis of a unit root.

**Remark 2.2** Theorem 2.1 can be also extended to cover the case of a deterministic trend. In particular, if the process under $H_0$ is generated by the equation $Y_t = X_t + a + bt$, where $(X_t)$ satisfies Assumption 2(i) and the regression equation

$$Y_t = \rho Y_{t-1} + a + bt + \sum_{j=1}^p a_{j,p} \Delta X_{t-j} + e_{t,p},$$

is fitted to the observed time series $Y_1, Y_2, \ldots, Y_n$, then the distribution of the least squares estimator $\hat{\rho}_n$ of $\rho$ is the same as the one given in Thorem 2.1 with the standard
Brownian motion $W(t)$ replaced by

$$\tilde{W}(t) = W(t) + (6t - 4) \int_0^1 W(s)ds - (12t - 6) \int_0^1 sW(s)ds.$$ 

### 2.3 Behavior Under the Alternative

Before studying the large-sample distribution of the least squares estimator $\hat{\rho}_n$ under the alternative of stationarity, we discuss an asymptotic collinearity problem that occurs when regression equation (1.1) is fitted to a time series stemming from a stationary process. This collinearity problem is essential for understanding the effects of choosing the truncation parameter $p$ on the power behavior of the test. The following proposition summarizes this behavior, and is of interest in its own right.

**Proposition 2.1** Let $\{W_t, t \in \mathbb{Z}\}$ be a zero mean, second order stationary process with autocovariance function $\gamma_W(h) = E(W_tW_{t+h})$ and spectral density $f_W$ satisfying $f_W(0) > 0$. Denote by $\mathcal{M}_{t,t-p} = \overline{\mathcal{P}}\{\Delta W_t, \Delta W_{t-1}, \cdots \Delta W_{t-p}\}$ the closed linear span generated by the differences $\Delta W_{t-j}, j = 0, 1, \ldots, p$ and by $\mathcal{P}_A(Y)$ the orthogonal projection of $Y$ onto the closed set $A$.

1. If $\gamma_W(h) \to 0$ as $h \to \infty$ then $E(W_t - \mathcal{P}_{\mathcal{M}_{t,t-p}}(W_t))^2 \to 0$ as $p \to \infty$.

2. If $\sum_{h=-\infty}^{\infty} |\gamma_W(h)| < \infty$ then $p \cdot E(W_t - \mathcal{P}_{\mathcal{M}_{t,t-p}}(W_t))^2 \to 2\pi f_W(0)$ as $p \to \infty$ where $f_W(\cdot)$ denotes the spectral density of $\{W_t\}$.

What the above proposition essentially says is that if $W_t$ is a second order stationary process, then $W_{t-1}$ can be expressed as a linear combination of its own differences $\Delta W_{t-j}, j = 1, 2, \ldots$. As alluded to in the Introduction, this proposition has serious consequences for the power behavior of the ADF-test under the alternative $H_1$. In particular, it implies a severe asymptotic collinearity problem that emerges when regression (1.1) is fitted to a stationary time series $X_1, X_2, \ldots, X_n$. 

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To elaborate, under the alternative of stationarity, the random variables $X_{t-1}$ and $\Delta X_{t-j}$, $j = 1, 2, \ldots$ appearing on the right hand side of (2.7) are perfectly collinear. Consequently, in fitting (the truncated) equation (1.1), the random variables $X_{t-1}$ and $\Delta X_{t-j}$, $j = 1, 2, \ldots, p$ which appear as regressors, become asymptotically collinear as the truncation parameter $p$ increases to infinity. Moreover, as part (ii) of Proposition 2.1 shows, the corresponding mean square prediction error $E(X_{t-1} - \mathcal{P}_{M_{t-1}, \ldots, \tilde{X}_{t-1}})^2$ is of order $O(1/p)$; this approximate collinearity and the resulting ill-conditioning may pose problems even for small values of $p$—as in our real data example of Section 3—especially when the covariance structure of the process is significant only at small lags.

As a final observation, note that the asymptotic collinearity problem occurs even if the underlying process $\{U_t\}$ is a finite, $p$-th order stationary autoregressive (AR) process and equation (1.1) is fitted to the observed time series using a truncation order $p_n$ that is allowed to increase (to infinity) as the sample size $n$ increases. Of course, if it is known that the process $\{U_t\}$ is a linear AR($p$) with known and finite $p$, then the truncation order $p_n$ would be held constant; this rather unrealistic situation seems to be the only case where the aforementioned collinearity problem disappears.

The following theorem establishes the limiting distribution of the least squares estimator $\hat{\rho}_n$ under the alternative $H_1$ in eq. (2.8).

**Theorem 2.2** Let Assumption 1 and Assumption 2(ii) be satisfied and suppose that $p = p_n \to \infty$ as $n \to \infty$ such that $p_n^4/\sqrt{n} \to 0$ and $\sqrt{n} \sum_{j=p+1}^{\infty} |a_j| \to 0$. Then, as $n \to \infty$,

(i) $\frac{n}{p} \text{Var}(\hat{\rho}_n) \to (1 - \rho)^2$, in probability, and

(ii) $\sqrt{n} \frac{p}{\sqrt{p}} (\hat{\rho}_n - \rho) \Rightarrow N(0, (1 - \rho)^2),$

where $\rho = \sum_{j=1}^{\infty} b_j$. 
Remark 2.3 [Rate of convergence of $\hat{\rho}_n$ under $H_1$] Notice that because in regression (1.1) we are interested in estimating the parameter $\rho$ only, we would expect, under the alternative of stationarity, that the estimator $\hat{\rho}_n$ would be $\sqrt{n}$-consistent. Surprisingly, this is not true. Delving into the matter deeper, it is apparent that the lower $\sqrt{n/p}$ convergence rate of $\hat{\rho}_n$ is due to the fact that estimating $\rho$ is tantamount to estimating the spectral density of $\{X_t\}$ at frequency zero. In fact, using $2\pi f_X(\lambda) = \sigma^2_\varepsilon / |1 - \sum_{j=0}^\infty b_j \exp \{i\lambda j\}|^2$, we get that $\rho = 1 - \sigma^2_\varepsilon / \sqrt{2\pi f_X(0)}$. This makes it clear that although $\rho$ appears to be a single parameter in the regression equation (1.1), estimating $\rho$ is essentially a nonparametric estimation problem. This behavior of the estimator $\hat{\rho}_n$ is regression (1.1) is different compared to the least squares estimator $\hat{a}_n$ in regression (2.9) considered by Phillips and Perron (1988) that is $\sqrt{n}$-consistent under $H_1$. The reason for the different convergence rates of the two estimators under $H_1$ lies in the fact that $\hat{\rho}_n$ in regression (1.1) estimates a function of the spectral density of $\{X_t\}$ at frequency zero, while $\hat{a}_n$ in regression (2.9) estimates the first order autocorrelation, see also Remark 2.1. Note, however, that the Phillips and Perron (1988) test suffers from a difficulty of its own in that its estimated critical value is typically not $\sqrt{n}$-consistent being itself a function of the underlying spectral density.

Condition $p_n^4 / \sqrt{n} \to 0$ in Theorem 2.2 seems unusual but it is not just an artifact of the proof; the order $p_n$ has to be small enough to avoid collinearity problems in the ADF regression, and condition $p_n^4 / \sqrt{n} \to 0$ ensures just that. It is of interest to see what happens when this condition is violated. From the proof of Theorem 2.2 it follows that under the alternative $H_1$ we have

$$\sqrt{n/p}(\hat{\rho} - \rho) = \frac{n^{-1/2} \sum_{t=p+1}^n V_{t-1,p} \varepsilon_t + O_P(p^{7/2}/\sqrt{n})}{\sigma^2_\varepsilon (1 - \rho)^{-2} + O_P(p^4/\sqrt{n})},$$

where $V_{t-1,p} = \sqrt{p}(X_{t-1} - \sum_{j=1}^p \delta_{j,p} \Delta X_{t-j})$, $\delta_{j,p}$, $j = 1, 2, \ldots, p$, are the coefficients of the best linear predictor of $X_{t-1}$ based on $\Delta X_{t-j}$, $j = 1, 2, \ldots, p$ and $n^{-1/2} \sum_{t=p+1}^n V_{t-1,p} \varepsilon_t =$
Corollary 2.1 Let Assumption 1 and Assumption 2(ii) be satisfied and suppose that 
p = p_n \to \infty \text{ as } n \to \infty \text{ such that } p_n/\sqrt{n} \to 0 \text{ but } p_n^4/\sqrt{n} \to \infty \text{ and } \sqrt{n} \sum_{j=p+1}^{\infty} |a_j| \to 0. \text{ Then, as } n \to \infty,
\sqrt{\frac{n}{p}}(\hat{\rho} - \rho) \to 0, \quad \text{in probability.}

Remark 2.4 [Comparison with previous results] Under the alternative of stationary ARMA processes, Lopez (1997) claimed that he has derived the asymptotic normal distribution of 
\sqrt{\frac{n}{p}} (\hat{\rho}_n - \rho) \text{ using the weaker condition } p^3/n \to 0 \text{ as } n \to \infty. \text{ Apart from the more restrictive ARMA process set-up, this statement is not justified. To elaborate, in order to derive the asymptotic distribution of } \sqrt{\frac{n}{p}} (\hat{\rho}_n - \rho) \text{ for this class of alternatives, Lopez (1997) proceeds by first replacing the ADF regression equation (1.1) by an autoregression containing only the levels of the } X_t \text{'s, i.e., instead of equation (1.1) he considers the autoregression equation } X_t = \phi_1,pX_{t-1} + \phi_2,pX_{t-2} + \ldots + \phi_{p,p}X_{t-p} + v_{t,p}; \text{ see equation (9) in Lopez (1997). Then, instead of the estimator } \hat{\rho}_n, \text{ he investigates the estimator } \hat{\phi}_n = \sum_{i=1}^{p} \hat{\phi}_{i,p}, \text{ where } \hat{\phi}_{i,p} \text{ is the least squares estimator of } \phi_{i,p} \text{ in the aforementioned autoregression containing only levels. Using results obtained by Berk (1974), the limiting distribution of } \hat{\phi}_n \text{ is then easily established allowing for the truncation lag } p \text{ to increase to infinity such that } p^3/n \to 0. \text{ However, the important step missing in this proof is the theoretical justification for the validity of this replacement in the regression considered. In fact, what one needs to show is that under the stated assumptions}

\sqrt{\frac{n}{p}}(\hat{\phi}_n - \hat{\rho}_n) \to 0, \quad \text{in probability.}

But Corollary 2.1 shows that the above is simply not true when } p_n^4/\sqrt{n} \to \infty \text{ and } p^3/n \to 0, \text{ e.g., if we let } p_n \sim n^a \text{ for some } a \in (1/8, 1/3). \text{ The non-equivalence of the}
two statistics $\hat{\phi}_n$ and $\hat{\rho}_n$ for large values of $p_n$ is further made apparent noting that the regression equation using only levels of the $X_t$’s does not suffer from the collinearity problems that are present in the regression equation (1.1) which also contains differences; see Proposition 2.1. It is exactly this potential collinearity under the alternative that requires the truncation lag $p$ in regression (1.1) to increase to infinity quite slow, i.e., satisfying $p_n^4/n \to 0$, in order to avoid ill-conditioning and obtain asymptotic normality of the estimator $\hat{\rho}_n$ as in our Theorem 2.2.

Remark 2.5 [Power and consistency of the ADF test] Theorem 2.2 allows for the following approximative expression for the power function of the ADF-test for fixed alternatives,

$$P_{H_1}(t_n < C_\alpha) \approx \Phi\left(\sqrt{\frac{n}{p}}\frac{\hat{\text{Std}}(\hat{\rho}_n)}{\sqrt{1 - \rho}}C_\alpha + \sqrt{\frac{n}{p}}\right) \approx \Phi\left(C_\alpha + \sqrt{\frac{n}{p}}\right),$$

(2.10)

where $C_\alpha$ denotes the upper $\alpha$-percentage point of the limiting distribution given in Theorem 2.1 and the second approximation follows since under $H_1$, $\hat{\text{Std}}(\hat{\rho}_n) = \sqrt{\frac{p}{n}}(1 - \rho)/\sqrt{n} + o_P(\sqrt{\frac{p}{n}})$. Therefore, as long as $p/n \to 0$, the ADF test will be consistent but with a rate of convergence smaller than the parametric rate $n^{1/2}$. Furthermore, as it is seen from (2.10), asymptotically the power of the test is not affected by the distance between $\rho$ and its value under the null hypothesis ($\rho = 1$) which is surprising in its own right. Finally, note that the dominand term in the power is $O(\sqrt{\frac{n}{p}})$ which is a decreasing function of $p$. This explains the empirically observed fact that increasing the truncation parameter $p$ in (1.1) leads to a drop of power of the ADF-test; see e.g. Figures 9.2 and 9.5 of Patterson (2011). Sections 3 and 4 in what follows give numerical evidence to this decrease of power as $p$ increases which can be quite pronounced even for small sample sizes.
3 A real data example

We now turn to a real data example that shows the power issues associated with the ADF test in practice. To this end, consider the dataset of Figure 1 that is extensively discussed in Example 1.2 of the well-known textbook by Shumway and Stoffer (2010). The data represent yearly average global temperatures deviations from 1880 to 2009 where the deviations are measured in degrees Celcius from the 1951-1980 average.

A familiar question is whether the data are trending and/or is temperature taking a ‘random walk’, i.e., does the temperature dataset have a unit root? Indeed, a linear trend can be readily noticed in the temperature dataset, and can help explain (at least in part) the strong autocorrelation characterizing the data pictured in Figure 2. However, we do not wish to focus on the Global Warming hypothesis here that would amount to checking the statistical significance of the linear trend. Rather, we want to test if there is a unit root process superimposed on the estimated trend whether the latter is negligible or not.

Using the tseries package in the R language, the P-values of the two aforementioned unit root tests were computed via fitting an equation that includes estimating a linear trend as in Remark 2.2. The P-value of the Phillips and Perron (1988) test was 0.01 while the P-value of the ADF test was 0.70. Obviously, this tremendous difference in the P-values raises serious concerns.

From the documentation of the tseries package it is made apparent that the adf.test function uses a default value for the order \( p \) given by the formula \( p = \lfloor (n - 1)^{1/3} \rfloor \); this implies a choice of \( p = 5 \) for our dataset where \( n = 130 \). We first note that this formula gives an acceptable rate for \( p \) under \( H_0 \) where it is just needed that \( p/\sqrt{n} \to 0 \); see Theorem 2.1. Nevertheless, a rate \( p \sim n^{1/3} \) seems too large under \( H_1 \); recall that Theorem 2.2 requires a rate satisfying \( p = o(n^{1/8}) \) in order to alleviate the collinearity and ill-conditioning of the ADF regression. The finite-sample simulations of Section 4 confirm the insights offered by the asymptotic Theorem 2.2 in that good
power seems to be associated with choosing $p$ quite small. Note that $n^{1/8} \approx 1.8$ when $n = 130$, so to have $p = o(n^{1/8})$ in our data example one would be well-advised to let $p = 1$.

The detrended data are shown in Figure 3, and their correlogram of Figure 4 does not show particularly strong dependence. Indeed, the estimated lag-1 autocorrelation is about 0.6 which does not give strong evidence for a unit root. Furthermore, the partial autocorrelation of the detrended data shown in Figure 5 suggests that a stationary AR(1) model might be quite appropriate—at least if one is ready to treat the lag-4 value as negligible for the purposes of parsimony. Running the ADF regression with the choice $p = 1$ actually results in a P-value just slightly under 0.01 that is in close agreement with the P-value of the Phillips and Perron (1988) test.

The above discussion helps underscore both the claimed loss of power associated with even a moderately large value of $p$ in the ADF test, as well as the need to scrutinize the choice of $p$ in practice as the implications can be quite severe. The simulations of Section 4 offer additional confirmation to these empirical findings.
Figure 2: Correlogram of yearly average global temperatures deviations.

Figure 3: The dataset of yearly average global temperatures deviations after removal of a linear trend.
Figure 4: Correlogram of the detrended dataset of yearly average global temperatures deviations.

Figure 5: Partial autocorrelation of the detrended dataset of yearly average global temperatures deviations.
4 A simulation experiment: good power vs. accurate size

In order to complement our asymptotic results, as well as shed some light on the unusual behavior of the ADF test on the real data example of Section 3, a numerical experiment was carried out. The data generating process was the ARMA (1,1) model:

\[ X_t - \phi X_{t-1} = Z_t + \theta Z_{t-1} \]

with \( Z_t \sim \text{i.i.d. } N(0, 1) \). For each of six different combinations of the ARMA parameters \( \phi \) and \( \theta \), and each sample size \( n \), 100,000 time series of length \( n \) were generated.

The simulation allowed us to compute empirical rejection probabilities of the ADF test at nominal level 0.05; these are shown in Tables 1–6. In terms of the chosen order, the formula \( p = na \) was used with a range of different values for \( a \). The last column from each Table corresponds to \( p \) chosen by minimizing the AIC criterion as detailed in McLeod et al. (2012, p. 687); the AIC minimization was carried out over orders \( p \) ranging from 1 to \( \sqrt{n} \).

Tables 1–2 concern the case where \( \{X_t\} \) has a unit root; hence, ideally all these entries should be close to the nominal level 0.05. The entry closest to 0.05 in each row is shown in boldface; for readability, the entries are rounded to three decimal points but in case of ties more decimal points were used to pinpoint the optimal \( p \). The standard error of each of those entries is quite small, about 0.0007, hence the entries give valuable and accurate information on the optimal \( p \).

Putting the results of Tables 1–2 together, it seems that the formula \( p = na \) with \( a \approx 0.3 \) works very well for accurate size of the ADF test. Interestingly, although AIC minimization gives generally good performance in the case of a positive MA parameter \( \theta \), in the case of a negative \( \theta \) it performs rather poorly except in the case of the huge sample with \( n = 1600 \).

Tables 3–6 concern the case where \( \{X_t\} \) does not have a unit root, i.e., under the
alternative; hence, the entries represent power, and should be as big as possible. As before, the biggest entries in each row are shown in boldface. The following conclusions can be drawn.

- Finite-sample powers are in general quite low; nevertheless, power tends to one as \( n \) increases. Therefore, the consistency of the test as claimed in Remark 2.5 is empirically verified.

- The \( p \) chosen to optimize size, i.e., \( p = n^{0.3} \), most definitely does \textit{not} optimize power.

- Even the formula \( p = n^a \) for some (other) constant \( a \) does not seem to work for optimal power; if anything, Tables 3–6 seem to suggest power advantages when using a formula of the type \( p = n^{a_n} \) with \( a_n \) being a decreasing function of \( n \).

- AIC minimization works reasonably well in terms of optimal power, in particular in the case of positive MA parameter \( \theta \).

\textbf{Remark 4.1} All in all, Tables 3–6 give empirical confirmation of the theoretical treatment of Section 2 pointing out the need to use small orders \( p \) in order to increase power and avoid collinearity. In addition, it is encouraging that AIC minimization works well to optimize power. Unfortunately, the no-free-lunch principle seems to apply here: implementing the ADF test using AIC minimization to chose the order may lead to a test that has wrong size, namely over-rejecting under the null; see e.g. Table 2.
Table 1. Entries represent empirical rejection probabilities of the ADF test using $p = n^a$ for different values of $a$; the last column corresponds to $p$ chosen by minimizing the AIC criterion. Data generating process was ARMA (1,1) with $(\phi, \theta) = (1,0.5)$.

Table 2. As in Table 1 with $(\phi, \theta) = (1,-0.5)$. 

23
\[(\phi, \theta) = (0.985, 0.5)\]

<table>
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<td>0.084</td>
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**Table 3.** As in Table 1 with \((\phi, \theta) = (0.985, 0.5)\).

\[(\phi, \theta) = (0.985, -0.5)\]

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**Table 4.** As in Table 1 with \((\phi, \theta) = (0.985, -0.5)\).
\[ (\phi, \theta) = (0.970, 0.5) \]

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**Table 5.** As in Table 1 with \( (\phi, \theta) = (0.970, 0.5) \).

\[ (\phi, \theta) = (0.970, -0.5) \]

<table>
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<th>0.09</th>
<th>0.13</th>
<th>0.17</th>
<th>0.21</th>
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<td><strong>0.171</strong></td>
<td>0.171</td>
<td>0.170</td>
<td>0.135</td>
<td>0.134</td>
<td>0.111</td>
<td>0.112</td>
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<td>0.356</td>
<td>0.334</td>
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</table>

**Table 6.** As in Table 1 with \( (\phi, \theta) = (0.970, -0.5) \).
5 Technical Proofs

**Proof of Theorem 2.1:** Using the notation $\Delta X_{t,p} = (\Delta X_t, \Delta X_{t-1}, \ldots, \Delta X_{t-p+1})'$, $\hat{\varepsilon}_{t,p} = X_t - \hat{\rho}_n X_{t-1} - \sum_{j=1}^{p} a_{j,p} \Delta X_{t-j}$ and $\hat{\varepsilon}_{t,p} = X_t - \hat{\rho}_n \hat{a}_{j,p} \Delta X_{t-j}$ with least squares estimators $\hat{\rho}_n$ and $\hat{a}_{j,p}$, $j = 1, 2, \ldots, p$, it is easily verified that

$$ t_n = (\hat{\rho}_n - 1)/\text{Std}(\hat{\rho}_n) = L_n R_n^{-1}/(\hat{\sigma}_n^2 R_n^{-1})^{1/2}, $$

where

$$ L_n = \sum_{t=p+1}^{n} X_{t-1} \hat{\varepsilon}_{t,p} - (\sum_{t=p+1}^{n} X_{t-1} \Delta X'_{t-1,p})(\sum_{t=p+1}^{n} \Delta X_{t-1,p} \Delta X'_{t-1,p})^{-1}(\sum_{t=p+1}^{n} \Delta X_{t,p} \hat{\varepsilon}_{t,p}), $$

$$ R_n = \sum_{t=p+1}^{n} X_{t-1}^2 - (\sum_{t=p+1}^{n} X_{t-1} \Delta X'_{t-1,p})(\sum_{t=p+1}^{n} \Delta X_{t-1,p} \Delta X'_{t-1,p})^{-1}(\sum_{t=p+1}^{n} \Delta X_{t,p} X_{t-1}) $$

and $\hat{\sigma}_n^2 = (n-p)^{-1} \sum_{t=p+1}^{n} \hat{\varepsilon}_{t,p}^2$ is the error variance estimator. Now, for $i, j \in \{1, 2, \ldots, p\}$ we have that, as $n \to \infty$, $n^{-1} \sum_{t=p+1}^{n} \Delta X_{t-i} \Delta X_{t-j} \to \gamma_U(i-j)$ in probability, and that, by the same arguments as in Berk (1974), p.493,

$$ \|n^{-1} \sum_{t=p+1}^{n} \Delta X_{t-1,p} \Delta X'_{t-1,p}\| = O_P(1), $$

where for a matrix $C$, the norm $\|C\| = \sup_{\|x\| \leq 1} \|Cx\|$ is used and $\|x\|$ denotes the Euclidean norm of the vector $x$. Furthermore,

$$ \frac{1}{np^{1/2}} \| \sum_{t=p+1}^{n} \Delta X_{t-1,p} X_{t-1} \| = (p^{-1} \sum_{j=1}^{p} (n^{-1} \sum_{t=p+1}^{n} \Delta X_{t-j} X_{t-1})^2)^{1/2} = O_P(1) $$

since $n^{-1} \sum_{t=p+1}^{n} \Delta X_{t-j} X_{t-1} = n^{-1} \sum_{t=p+1}^{n} \sum_{j=1}^{n} U_{t-j} U_t = O_P(1)$. Finally, since

$$ n^{-1/2} \sum_{t=p+1}^{n} \Delta X_{t-j} \hat{\varepsilon}_{t,p} = n^{-1/2} \sum_{t=p+1}^{n} U_{t-j} \hat{\varepsilon}_{t,p} = O_P(1) $$

we get that

$$ \sqrt{p} n^{-1/2} \| \sum_{t=p+1}^{n} \Delta X_{t-1,p} \hat{\varepsilon}_{t,p} \| = \frac{p}{\sqrt{n}} (p^{-1} \sum_{j=1}^{p} (1/\sqrt{n} \sum_{t=p+1}^{n} \Delta X_{t-j} \hat{\varepsilon}_{t,p})^2)^{1/2} = \frac{p}{\sqrt{n}} O_P(1) \to 0 $$

(5.15)
as \( n \to \infty \). Now, equations (5.13) to (5.15) implies that
\[
n^{-1}L_n = n^{-1} \sum_{t=p+1}^{n} X_{t-1}\varepsilon_{t,p} + o_P(1).
\]
Furthermore, because
\[
n^{-1} \sum_{t=p+1}^{n} X_{t-1}(\varepsilon_{t,p} - \varepsilon_t) = \sum_{j=1}^{p} (a_{j,p} - a_j)n^{-1} \sum_{t=p+1}^{n} X_{t-1}U_{t-j} + \sum_{j=p+1}^{\infty} a_j n^{-1} \sum_{t=p+1}^{n} X_{t-1}U_{t-j},
\]
if follows using \( n^{-1} \sum_{t=p+1}^{n} X_{t-1}U_{t-j} = O_P(1) \) and Baxter’s inequality, cf. Lemma 2.2 of Kreiss et al. (2011), that
\[
|n^{-1} \sum_{t=p+1}^{n} X_{t-1}(\varepsilon_{t,p} - \varepsilon_t)| \leq O_P(\sum_{j=p+1}^{\infty} |a_j|) \to 0,
\]
as \( p \to \infty \). Thus,
\[
n^{-1}L_n = n^{-1} \sum_{t=p+1}^{n} X_{t-1}\varepsilon_t + o_P(1). \tag{5.16}
\]
Similarly, using (5.13) and (5.14) we obtain that
\[
n^{-2}R_n = n^{-2} \sum_{t=p+1}^{n} X_{t-1}^2 + o_P(1). \tag{5.17}
\]

Now, as in the proof of Theorem 3.1 in Phillips (1987) and using the invariance principle for the partial sum process \( S_{[nr]} = n^{-1/2} \sum_{j=1}^{[nr]} \varepsilon_j \) of zero mean weakly dependent random variables satisfying Assumption 1(ii), established in Theorem 1 of Wu and Min (2005), we get that
\[
n^{-1}L_n \Rightarrow \sigma^2_{\varepsilon} \int_{0}^{1} W(t)dW(t), \quad \text{and} \quad n^{-2}R_n \Rightarrow \sigma^2_{\varepsilon} \int_{0}^{1} W^2(t)dt.
\]

\[ \square \]

**Proof of Proposition 2.1:** Let \( \delta_{j,p}, \ j = 0, 1, \ldots, p \) be the coefficients of \( \Delta X_{t-j} \) in the best linear prediction of \( X_t \) based on \( \Delta X_{t-j} \) and define \( l_{j,p} = (1 - j/p) \), \( j = 0, 1, \ldots, p \).

We have
\[
E(X_t - P_{M_{t-p}} X_t)^2 = E(X_t - \sum_{j=0}^{p} l_{j,p} \Delta X_{t-j})^2 + E(\sum_{j=0}^{p} (l_{j,p} - \delta_{j,p}) \Delta X_{t-j})^2
\]
\[
- 2E(\sum_{j=0}^{p} (l_{j,p} - \delta_{j,p}) \Delta X_{t-j})(X_t - \sum_{j=0}^{p} l_{j,p} \Delta X_{t-j}).
\]

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Note first that \( E(X_t - \sum_{j=0}^p l_{j,p} \Delta X_{t-j})^2 = p^{-1}[\gamma(0) + 2 \sum_{s=1}^{p-1} (1 - s/p) \gamma(s)] \), which converges to zero if \( \gamma(h) \to 0 \). Furthermore, if \( \sum_{h=-\infty}^{\infty} |\gamma(h)| < \infty \) then \( p \cdot E(X_t - \sum_{j=0}^p l_{j,p} \Delta X_{t-j})^2 \to \sum_{h=-\infty}^{\infty} \gamma(h) = 2\pi f_{X_1}(0) \) by the dominate convergence theorem.

Now, let \( X_t(p) = (X_t, X_{t-1}, \ldots, X_{t-p})', \Gamma_{p+1} = E(X_t(p)X_t(p)') \) and define the \((p+1)\)-dimensional vectors \( \tilde{\delta}(p) = ((1 - \delta_{0,p}), (\delta_{0,p} - \delta_{1,p}), \ldots, (\delta_{p-1,p} - \delta_{p,p}), \delta_{p,p}) \) and \( \tilde{l}(p) = (0, 1/p, 1/p, \ldots, 1/p, 0)' \). Then the following upper bound is valid,

\[
E(\sum_{j=0}^p (l_{j,p} - \delta_{j,p}) \Delta X_{t-j})^2 = (\tilde{l}(p) - \tilde{\delta}(p))' \Gamma_{p+1} (\tilde{l}(p) - \tilde{\delta}(p)) \leq \max_{\lambda \in [0, \pi]} f_{X_1}(\lambda) \| (\tilde{l}(p) - \tilde{\delta}(p)) \|^2 \leq \max_{\lambda \in [0, \pi]} f_{X_1}(\lambda) \left( 2\| \tilde{l}(p) \|^2 + 2\| \tilde{\delta}(p) \|^2 \right).
\]

It is easily seen that \( \| \tilde{l}(p) \|^2 = O(p^{-1}) \to 0 \). Furthermore, using the following lower bound for the mean square prediction error

\[
E(X_t - P_{M_t,t-p} X_t)^2 = \int_{-\pi}^{\pi} | \sum_{j=0}^p \delta_{j,p} e^{-ij\lambda} |^2 f_{X_1}(\lambda) d\lambda \geq \inf_{\lambda \in [0, \pi]} f_{X_1}(\lambda) \| \tilde{\delta}(p) \|^2,
\]

we get that \( \| \tilde{\delta}(p) \|^2 \to 0 \) as \( p \to \infty \) from which it follows that \( E(\sum_{j=0}^p (l_{j,p} - \delta_{j,p}) \Delta X_{t-j})^2 \to 0 \) as \( p \to 0 \). Finally, by the above results and Cauchy-Schwarz’s inequality, it follows that \( |E(\sum_{j=0}^p (l_{j,p} - \delta_{j,p}) \Delta X_{t-j})(X_t - \sum_{j=0}^p l_{j,p} \Delta X_{t-j})| \to 0 \) which concludes the proof.

\[\Box\]

**Proof of Theorem 2.2:** Note that \( \sqrt{n/p}(\hat{\gamma}_n - \gamma) = \sqrt{n/p} L_n R_n^{-1} \) where \( L_n \) and \( R_n \) are defined in (5.11) and (5.12). Let \( \hat{\gamma}_0 = (n-p)^{-1} \sum_{t=p+1}^n X_{t-1}^2 \),

\[
\hat{d}_p = \left( \frac{1}{n-p} \sum_{t=p+1}^n \Delta X_{t-i} X_{t-1}, i = 1, 2, \ldots, p \right)', \quad \text{and} \quad \hat{C}_p = \left( \frac{1}{n-p} \sum_{t=p+1}^n \Delta X_{t-i} \Delta X_{t-j} \right)_{i,j=1,2,\ldots,p}.
\]

We have that

\[
n^{-1} R_n = \hat{\gamma}_0 - \hat{d}_p \hat{C}_p^{-1} \hat{d}_p = \gamma_0 - d_p' C_p^{-1} d_p + O_P(p^3/n^{1/2}), \quad (5.18)
\]
where $d'_p = (E(X_{t-1} \Delta X_{t-j}), j = 1, 2, \ldots, p)$ and $C_p = E(\Delta X_{t-1,p} \Delta X'_{t-1,p})$. Notice that the $O_P(p^2/n^{1/2})$ term in (5.18) appears because using the notation $\tau_p^2 = \gamma_0 - d'_p C_p^{-1} d_p$ and $\tilde{\tau}_p^2 = \gamma_0 - \hat{d}_p \hat{C}_p^{-1} \hat{d}_p$, we have that

$$|\tilde{\tau}_p^2 - \tau_p^2| \leq |\gamma_0 - \gamma_0| + \|\hat{\delta}_p - \delta_p\| \|d_p\| + \|\hat{d}_p - d_p\| \|\hat{\delta}_p\|$$

where $\hat{\delta}_p = \hat{C}_p^{-1} \hat{d}_p$ and $\delta_p = C_p^{-1} d_p$. Now, $\|\hat{d}_p - d_p\| = O_P(p^{1/2}/n^{1/2})$ and

$$p^{-1/2} \|\hat{d}_p\| = \{p^{-1} \sum_{j=1}^{p} ((n - p)^{-1} \sum_{t=p+1}^{n} \Delta X_{t-j} X_{t-1})^2 \}^{1/2} = O_P(1).$$

Furthermore,

$$\|\hat{\delta}_p - \delta_p\| = O_P(p^{5/2}/n^{1/2}), \tag{5.19}$$

and

$$p^{-1} \|\hat{\delta}_p\| = o_P(1), \tag{5.20}$$

which implies that $|\tilde{\tau}_p^2 - \tau_p^2| = O_P(p^3/n^{1/2})$. We show that (5.19) and (5.20) are true.

To see (5.19) notice first that

$$\hat{\delta}_p - \delta_p = \hat{C}_p^{-1} \{(n - p)^{-1} \sum_{t=p+1}^{n} \Delta X_{t-1,p} u_{t-1,p}\},$$

where $u_{t-1,p} = X_{t-1} - \sum_{j=1}^{p} \delta_{j,p} \Delta X_{t-j}$ and $\delta_p = (\delta_{1,p}, \delta_{2,p}, \ldots, \delta_{p,p})$ are the coefficients of the best linear predictor of $X_{t-1}$ based on $\Delta X_{t-j}, j = 1, 2, \ldots, p$. Now,

$$\|\hat{\delta}_p - \delta_p\| \leq \|\hat{C}_p^{-1}\| \|n - p\|^{-1} \sum_{t=p+1}^{n} \Delta X_{t-1,p} u_{t-1,p} = O_P(p^{5/2}/n^{1/2}),$$

since $\|n - p\|^{-1} \sum_{t=p+1}^{n} \Delta X_{t-1,p} u_{t-1,p} = O_P(p^{1/2}/n^{1/2})$, and

$$\|\hat{C}_p^{-1}\| = O_P(p^2). \tag{5.21}$$

To see why (5.21) is true notice that for every $p \in \mathbb{N}$ the matrix $C_p$ is positive definite, $\|C_p^{-1}\|$ is the reciprocal of the minimal eigenvalue of $C_p$. Notice that the spectral density
where $\tilde{\lambda}_{\text{min}}$ denotes the minimal eigenvalue of the $p \times p$ covariance matrix of the process with spectral density $(2\pi)^{-1}|1-e^{i\lambda}|^2 = (2\pi)^{-1}2(1-\cos(\lambda))$, i.e., of the noninvertible MA(1) process $Y_t = \varepsilon_t - \varepsilon_{t-1}$. For this process the eigenvalues of the $p$-dimensional correlation matrix are given by $\tilde{\lambda}_k = 2(1-\cos((k\pi)/(p+1)))$, $k = 1, 2, \ldots, p$. Thus, $\|C_p^{-1}\| \leq 1/K(1-\cos(\pi/(p+1))) = O(p^2)$, where the last equality follows because $[1-\cos(\pi/(p+1))] \sim p^{-2}$. Note that $\|C_p^{-1}\| \to \infty$ as $p \to \infty$ since the minimal eigenvalue of $C_p$ approaches zero as $p \to \infty$. Furthermore, it is easily seen that $\|\hat{C}_p - C_p\| = O(p/\sqrt{n})$ from which we get using $\|C_p^{-1}\| = O(p^2)$ that

$$
\|\hat{C}_p^{-1} - C_p^{-1}\| \leq \frac{\|C_p^{-1}\|^2\|\hat{C}_p - C_p\|}{1 - \|\hat{C}_p - C_p\|\|C_p^{-1}\|} = O(p^{5/2}/n^{1/2}).
$$

Thus,

$$
\|\hat{C}_p^{-1}\| \leq \|C_p^{-1}\| + \|\hat{C}_p^{-1} - C_p^{-1}\| = O(p^2 + p^{5/2}/n^{1/2}).
$$

To see (5.20) notice that

$$
\|\hat{\delta}_p\| \leq \|\delta_p\| + \|\hat{\delta}_p - \delta_p\|.
$$

Now, $\|\hat{\delta}_p - \delta_p\| = O(p^{5/2}/n^{1/2})$. Furthermore, for $l_p = (l_{1,p}, l_{2,p}, \ldots, l_{p,p})'$, $l_{j,p} = (1 - j/p)$, $j = 1, 2, \ldots, p$ we have $\|\delta_p\| \leq \|l_p\| + \|\delta_p - l_p\|$. Now, $\|l_p\| = O(\sqrt{p})$ while $\|l_p - \delta_p\| = o(p)$ which follows because

$$
(l_p - \delta_p)'C_p(l_p - \delta_p) \geq \lambda_{\text{min}}\|l_p - \delta_p\|^2 \geq \inf_{\lambda \in [0,\pi]} f_X(\lambda)(1 - \cos(\pi/(p+1))\|l_p - \delta_p\|^2 \\
\sim Kp^{-2}\|l_p - \delta_p\|^2 \geq 0,
$$

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and as the proof of Proposition 2.1 shows $(l_p - d_p)'C_p(l_p - d_p) = E(\sum_{j=0}^{p}(l_{j,p} - \delta_{j,p})\Delta X_{t-j})^2 \to 0$, as $p \to \infty$.

This concludes the proof of assertion (5.18).

Now,

$$p\tau_p^2 \to 2\pi f_{X_t}(0) = \sigma^2(1 - \rho)^{-2},$$

by Proposition 2.1. Thus

$$\frac{p}{n}R_n = \sigma^2(1 - \rho)^{-2} + O_P(p^4/n^{1/2}).$$

(5.23)

Let $\hat{V}_{t-1, p} = \sqrt{p}(X_{t-1} - \hat{d}_p \hat{C}_p^{-1} \Delta X_{t-1, p})$ and $V_{t-1, p} = \sqrt{p}(X_{t-1} - d_p C_p^{-1} \Delta X_{t-1, p})$. Then,

$$\sqrt{p}L_n = \frac{1}{\sqrt{n}} \sum_{t=p+1}^{n} \hat{V}_{t-1, p} \varepsilon_{t, p}$$

$$= \frac{1}{\sqrt{n}} \sum_{t=p+1}^{n} V_{t-1, p} \varepsilon_{t} + \frac{1}{\sqrt{n}} \sum_{t=p+1}^{n} (\hat{V}_{t-1, p} - V_{t-1, p}) \varepsilon_{t}$$

$$+ \frac{1}{\sqrt{n}} \sum_{t=p+1}^{n} \hat{V}_{t-1, p} (\varepsilon_{t, p} - \varepsilon_{t})$$

$$= \frac{1}{\sqrt{n}} \sum_{t=p+1}^{n} V_{t-1, p} \varepsilon_{t} + L_{1,n} + L_{2,n},$$

with an obvious notation for the remainder terms $L_{1,n}$ and $L_{2,n}$. For these terms we have

$$|L_{1,n}| \leq \sqrt{p} ||\hat{d}_p - \delta_p|| n^{-1/2} \sum_{t=p+1}^{n} \Delta X_{t-1, p} \varepsilon_{t} ||$$

$$= O_P(p^{7/2}/n^{1/2}) \to 0 \text{ as } n \to \infty. $$

and

$$|L_{2,n}| \leq (n^{-1} \sum_{t=p+1}^{n} \hat{V}_{t-1, p}^2)^{1/2} (\sum_{t=p+1}^{n} (\varepsilon_{t, p} - \varepsilon_{t})^2)^{1/2}$$

$$= O_P(\sqrt{n} \sum_{j=p+1}^{n} |a_j|).$$
Notice that the last equality above follows since under the assumptions made,

\[ n^{-1} \sum_{t=p+1}^{n} \hat{V}_{t-1,p}^2 = n^{-1} \sum_{t=p+1}^{n} p(X_{t-1} - \delta'_p \Delta X_{t-1,p})^2 + o_P(1) = O_P(1), \]

and by Baxter’s inequality, see Kreiss et al. (2011), Lemma 2.2,

\[ E(\sum_{t=p+1}^{n} (\varepsilon_{t,p} - \varepsilon_t)^2) = O(\sqrt{n} \sum_{j=p+1}^{\infty} |a_j|). \]

The proof of the theorem is then concluded since under the assumptions made and by a central limit theorem for martingale differences, see Theorem 1 of Brown(1971),

\[ \frac{1}{\sqrt{n}} \sum_{t=p+1}^{n} V_{t-1,p} \varepsilon_t \Rightarrow N(0, \sigma^4\varepsilon(1 - \rho)^{-2}). \quad (5.24) \]

References


