# Time Series Analysis and Its Applications 

With R Examples
EZ Edition

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With R Examples

## EZ Green Edition

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## Preface

The EZ version can be used for an introductory time series course where the prerequisite is an understanding of linear regression analysis, although it assumes some basic probability skills (expectation) and general high school math skills (trigonometry, complex numbers, polynomials, calculus, and so on):

- A calculus based introduction to probability course is an essential co-requisite for this course. Readers should be familiar with most of the content of basic probability facts (a pdf file), which is a short introduction to the necessary material.
- For readers who are a little rusty on the basic math skills, the WikiBook on K-12 mathematics http://en.wikibooks.org/wiki/Subject:K-12_mathematics may be a useful resource; in particular, we mention the book covering calculus. We occasionally use matrix notation. For readers lacking this skill, see the high school page on matrices. For Chapter 4, this primer on complex numbers may be helpful.

The EZ version follows the basic outline of the blue version, but with the addition of examples to take the place of entire sections. Appendix R has information regarding the use of the R package for the text astsa. The code listing on the website for the yellow text does not match the EZ version.

Two stars $(* *)$ indicate that skills obtained in a course on basic mathematical statistics are recommended and these parts may be skipped. Chapter 5 is a $* *-$ ed chapter on some advanced time domain topics. In this version, the section on regression with autocorrelated errors is in Chapter 3. In Chapter 3, there are two estimation sections, one that uses only least squares and method of moments, and a **-ed section on MLE that comes after the section on forecasting.

Internal links are dark red, external links are magenta, R code is in blue, output is purple and comments are \# green.

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## Chapter 1

## Time Series Characteristics

### 1.1 Introduction

The analysis of experimental data that have been observed at different points in time leads to new and unique problems in statistical modeling and inference. The obvious correlation introduced by the sampling of adjacent points in time can severely restrict the applicability of the many conventional statistical methods traditionally dependent on the assumption that these adjacent observations are independent and identically distributed. The systematic approach by which one goes about answering the mathematical and statistical questions posed by these time correlations is commonly referred to as time series analysis.

Historically, time series methods were applied to problems in the physical and environmental sciences. This fact accounts for the basic engineering flavor permeating the language of time series analysis. In our view, the first step in any time series investigation always involves careful scrutiny of the recorded data plotted over time. Before looking more closely at the particular statistical methods, it is appropriate to mention that two separate, but not necessarily mutually exclusive, approaches to time series analysis exist, commonly identified as the time domain approach (Chapter 3) and the frequency domain approach (Chapter 4).

### 1.2 Some Time Series Data

The following examples illustrate some of the common kinds of time series data as well as some of the statistical questions that might be asked about such data.

## Example 1.1 Johnson \& Johnson Quarterly Earnings

Figure 1.1 shows quarterly earnings per share for the U.S. company Johnson \& Johnson. There are 84 quarters ( 21 years) measured from the first quarter of 1960


Fig. 1.1. Johnson \& Johnson quarterly earnings per share, 1960-I to 1980-IV.


Fig. 1.2. Left: Quarterly value of initial deposits of $\$ 75, \$ 100, \$ 125$, and $\$ 150$ over 15 years, with a quarterly growth rate of $5 \% ; x_{t}=(1+.05) x_{t-1}$. Right: Logs of the quarterly values; $\log \left(x_{t}\right)=\log (1+.05)+\log \left(x_{t-1}\right)$. When marked in terms of quarters, Figure 1.2 looks like Figure 1.1.
to the last quarter of 1980. Modeling such series begins by observing the primary patterns in the time history. In this case, note the increasing underlying trend and variability, and a somewhat regular oscillation superimposed on the trend that seems to repeat over quarters. Methods for analyzing data such as these are explored in Chapter 2 (see Problem 2.1) using regression techniques. Also, compare Figure 1.1 with Figure 1.2. To use package astsa, and then plot the data for this example using R, type the following (try plotting the logged the data yourself).
require(astsa) \# ** SEE FOOTNOTE
plot(jj, type="o", ylab="Quarterly Earnings per Share")
plot $(\log (j j))$ \# not shown

## Example 1.2 Global Warming

Consider the global temperature series record shown in Figure 1.3. The data are the global mean land-ocean temperature index from 1880 to 2009, with the base period 1951-1980. In particular, the data are deviations, measured in degrees centigrade, from the 1951-1980 average, and are an update of Hansen et al. (2006). We note an apparent upward trend in the series during the latter part of the twentieth century

[^0]

Fig. 1.3. Yearly average global temperature deviations (1880-2009) in ${ }^{\circ} \mathrm{C}$.
that has been used as an argument for the global warming hypothesis. Note also the leveling off at about 1935 and then another rather sharp upward trend at about 1970. The question of interest for global warming proponents and opponents is whether the overall trend is natural or whether it is caused by some human-induced interface. The R code for this example is similar to the code in Example 1.1:
plot(gtemp, type="o", ylab="Global Temperature Deviations")

## Example 1.3 New York Stock Exchange

As an example of financial time series data, Figure 1.4 shows the daily returns (or percent change) of the New York Stock Exchange (NYSE) from February 2, 1984 to December 31, 1991. It is easy to spot the crash of October 19, 1987 in the figure. The data shown in Figure 1.4 are typical of return data. The mean of the series appears to be stable with an average return of approximately zero, however, the volatility (or variability) of data changes over time. In fact, the data show volatility clustering; that is, highly volatile periods tend to be clustered together. A problem in the analysis of these type of financial data is to forecast the volatility of future returns. For example, GARCH models have been developed to handle these problems. The R code for this example is similar to the previous examples:
plot(nyse, ylab="NYSE Returns")


Fig. 1.4. Daily weighted market returns from February 2, 1984 to December 31, 1991. The crash of October 19, 1987 occurs at $t=938$.


Fig. 1.5. Monthly SOI and Recruitment (estimated new fish), 1950-1987.

## Example 1.4 El Niño and Fish Population

We may also be interested in analyzing several time series at once. Figure 1.5 shows monthly values of an environmental series called the Southern Oscillation Index (SOI) and associated Recruitment (an index of the number of new fish). Both series are for a period of 453 months ranging over the years 1950-1987. SOI measures changes in air pressure related to sea surface temperatures in the central Pacific Ocean. The central Pacific warms every three to seven years due to the El Niño effect, which has been blamed, in particular, for the 1997 floods in the midwestern portions of the United States. Both series in Figure 1.5 tend to exhibit repetitive behavior, with regularly repeating cycles that are easily visible. This periodic behavior is of interest because underlying processes of interest may be regular and the rate or frequency of oscillation characterizing the behavior of the underlying series would help to identify them. The Recruitment series also shows several kinds of oscillations. The study of the kinds of cycles and their strengths is the subject of Chapter 4. The two series also tend to be somewhat related; it is easy to imagine that somehow the fish population is dependent on the SOI.

The following R code will reproduce Figure 1.5:

```
par(mfrow = c(2,1)) # set up the graphics
plot(soi, ylab="", xlab="", main="Southern Oscillation Index")
plot(rec, ylab="", xlab="", main="Recruitment")
```


## Example 1.5 fMRI Imaging

A fundamental problem in classical statistics occurs when we are given a collection of independent series or vectors of series, generated under varying experimental conditions or treatment configurations. Such a set of series is shown in Figure 1.6,


Fig. 1.6. fMRI data from various locations in the cortex, thalamus, and cerebellum; $n=128$ points, one observation taken every 2 seconds.


#### Abstract

where we observe data collected from various locations in the brain via functional magnetic resonance imaging (fMRI). In this example, a stimulus was applied for 32 seconds and then stopped for 32 seconds; thus, the signal period is 64 seconds. The sampling rate was one observation every 2 seconds for 256 seconds ( $n=128$ ). The series shown in Figure 1.6 are consecutive measures of blood oxygenation-level dependent (bоцд) signal intensity, which measures areas of activation in the brain. Notice that the periodicities appear strongly in the motor cortex series and less strongly in the thalamus and cerebellum. The fact that one has series from different areas of the brain suggests testing whether the areas are responding differently to the brush stimulus. Use the following R commands to plot the data: ```par(mfrow=c(2,1), mar=c (3,2,1,0)+.5, mgp=c (1.6,.6,0))``` ts.plot(fmri1[,2:5], col=1:4, ylab="BOLD", xlab="", main="Cortex") ts.plot(fmri1[,6:9], col=1:4, ylab="BOLD", xlab="", main="Thalam \& Cereb") mtext("Time (1 pt = 2 sec$) "$, side=1, line=2)


### 1.3 Time Series Models

The primary objective of time series analysis is to develop mathematical models that provide plausible descriptions for sample data, like that encountered in the previous section.

The fundamental visual characteristic distinguishing the different series shown in Example 1.1 - Example 1.5 is their differing degrees of smoothness. One possible explanation for this smoothness is that adjacent points in time are correlated, so
the value of the series at time $t$, say, $x_{t}$, depends in some way on the past values $x_{t-1}, x_{t-2}, \ldots$. This model expresses a fundamental way in which we might think about generating realistic-looking time series.

## Example 1.6 White Noise (3 flavors)

A simple kind of generated series might be a collection of uncorrelated random variables, $w_{t}$, with mean 0 and finite variance $\sigma_{w}^{2}$. The time series generated from uncorrelated variables is used as a model for noise in engineering applications where it is called white noise; we shall sometimes denote this process as $\boldsymbol{w}_{\boldsymbol{t}} \sim \boldsymbol{w n}\left(\mathbf{0}, \boldsymbol{\sigma}_{\boldsymbol{w}}^{2}\right)$. The designation white originates from the analogy with white light and indicates that all possible periodic oscillations are present with equal strength.

We will, at times, also require the noise to be independent and identically distributed (iid) random variables with mean 0 and variance $\sigma_{w}^{2}$. We will distinguish this by saying white independent noise, or by writing $\boldsymbol{w}_{\boldsymbol{t}} \sim \operatorname{iid}\left(\mathbf{0}, \boldsymbol{\sigma}_{\boldsymbol{w}}^{2}\right)$. A particularly useful white noise series is Gaussian white noise, wherein the $w_{t}$ are independent normal random variables, with mean 0 and variance $\sigma_{w}^{2}$; or more succinctly, $\boldsymbol{w}_{\boldsymbol{t}} \sim \operatorname{iid} \mathbf{N}\left(\mathbf{0}, \boldsymbol{\sigma}_{\boldsymbol{w}}^{\mathbf{2}}\right)$. Figure 1.7 shows in the upper panel a collection of 500 such random variables, with $\sigma_{w}^{2}=1$, plotted in the order in which they were drawn. The resulting series bears a slight resemblance to portions of the NYSE returns in Figure 1.4. The plot tends to show visually a mixture of many different kinds of oscillations in the white noise series.

If the stochastic behavior of all time series could be explained in terms of the white noise model, classical statistical methods would suffice. Two ways of introducing serial correlation and more smoothness into time series models are given in Example 1.7 and Example 1.8.

## Example 1.7 Moving Averages and Filtering

We might replace the white noise series $w_{t}$ by a moving average that smooths the series. For example, consider replacing $w_{t}$ in Example 1.6 by an average of its current value and its immediate neighbors in the past and future. That is, let

$$
\begin{equation*}
v_{t}=\frac{1}{3}\left(w_{t-1}+w_{t}+w_{t+1}\right) \tag{1.1}
\end{equation*}
$$

which leads to the series shown in the lower panel of Figure 1.7. Inspecting the series shows a smoother version of the first series, reflecting the fact that the slower oscillations are more apparent and some of the faster oscillations are taken out. We begin to notice a similarity to some of the non-cyclic fMRI series in Figure 1.6.

To reproduce Figure 1.7 in R use the following commands. A linear combination of values in a time series such as in (1.1) is referred to, generically, as a filtered series; hence the command filter.

```
w = rnorm(500,0,1) # 500 N(0,1) variates
v = filter(w, sides=2, rep(1/3,3)) # moving average
par(mfrow=c(2,1))
plot.ts(w, main="white noise")
plot.ts(v, ylim=c(-3,3), main="moving average")
```



Fig. 1.7. Gaussian white noise series (top) and three-point moving average of the Gaussian white noise series (bottom).

The SOI and Recruitment series in Figure 1.5, as well as some of the MRI series in Figure 1.6, differ from the moving average series because one particular kind of oscillatory behavior seems to predominate, producing a sinusoidal type of behavior. A number of methods exist for generating series with this quasi-periodic behavior; we illustrate a popular one based on the autoregressive model considered in Chapter 3.

## Example 1.8 Autoregressions

Suppose we consider the white noise series $w_{t}$ of Example 1.6 as input and calculate the output using the second-order equation

$$
\begin{equation*}
x_{t}=x_{t-1}-.9 x_{t-2}+w_{t} \tag{1.2}
\end{equation*}
$$

successively for $t=1,2, \ldots, 500$. Equation (1.2) represents a regression or prediction of the current value $x_{t}$ of a time series as a function of the past two values of the series, and, hence, the term autoregression is suggested for this model. A problem with startup values exists here because (1.2) also depends on the initial conditions $x_{0}$ and $x_{-1}$, but, for now, we assume that we are given these values and generate the succeeding values by substituting into (1.2). That is, given $w_{1}, w_{2}, \ldots, w_{500}$, and $x_{0}, x_{-1}$, we start with $x_{1}=x_{0}-.9 x_{-1}+w_{1}$, then recursively compute $x_{2}=x_{1}-.9 x_{0}+w_{2}$, then $x_{3}=x_{2}-.9 x_{1}+w_{3}$, and so on. The resulting output series is shown in Figure 1.8, and we note the periodic behavior of the series, which is similar to that displayed by the SOI and Recruitment in Figure 1.5 and some fMRI series in Figure 1.6. The autoregressive model above and its generalizations can be used as an underlying model for many observed series and will be studied in detail in Chapter 3.


Fig. 1.8. Autoregressive series generated from model (1.2).

One way to simulate and plot data from the model (1.2) in R is to use the following commands (another way is to use arima.sim). The initial conditions are set equal to zero, so we let the filter run an extra 50 values to avoid startup problems.
$\mathrm{w}=\operatorname{rnorm}(550,0,1)$ \# 50 extra to avoid startup problems
$\mathrm{x}=$ filter(w, filter=c(1,-.9), method="recursive") [-(1:50)]
plot.ts(x, main="autoregression")

## Example 1.9 Random Walk with Drift

A model for analyzing trend such as seen in the global temperature data in Figure 1.3, is the random walk with drift model given by

$$
\begin{equation*}
x_{t}=\delta+x_{t-1}+w_{t} \tag{1.3}
\end{equation*}
$$

for $t=1,2, \ldots$, with initial condition $x_{0}=0$, and where $w_{t}$ is white noise. The constant $\delta$ is called the drift, and when $\delta=0$, the model is called simply a random walk because the value of the time series at time $t$ is the value of the series at time $t-1$ plus a completely random movement determined by $w_{t}$. Note that we may rewrite (1.3) as a cumulative sum of white noise variates. That is,

$$
\begin{equation*}
x_{t}=\delta t+\sum_{j=1}^{t} w_{j} \tag{1.4}
\end{equation*}
$$

for $t=1,2, \ldots$; either use induction, or plug (1.4) into (1.3) to verify this statement. Figure 1.9 shows 200 observations generated from the model with $\delta=0$ and .2 , and with $\sigma_{w}=1$. For comparison, we also superimposed the straight lines $\delta t$ on the graph.

To reproduce Figure 1.9 in R use the following code (notice the use of multiple commands per line using a semicolon).

```
set.seed(154) # so you can reproduce the results
w = rnorm(200,0,1); x = cumsum(w) # two commands in one line
wd = w +.2; xd = cumsum(wd)
plot.ts(xd, ylim=c(-5,55), main="random walk")
lines(x); lines(.2*(1:200), lty="dashed")
```



Fig. 1.9. Random walk, $\sigma_{w}=1$, with drift $\delta=.2$ (upper jagged line), without drift, $\delta=0$ (lower jagged line), and dashed lines showing the drifts.

## Example 1.10 Signal in Noise

Many realistic models for generating time series assume an underlying signal with some consistent periodic variation, contaminated by adding a random noise. For example, it is easy to detect the regular cycle fMRI series displayed on the top of Figure 1.6. Consider the model

$$
\begin{equation*}
x_{t}=2 \cos \left(2 \pi \frac{t+15}{50}\right)+w_{t} \tag{1.5}
\end{equation*}
$$

for $t=1,2, \ldots, 500$, where the first term is regarded as the signal, shown in the upper panel of Figure 1.10. We note that a sinusoidal waveform can be written as

$$
\begin{equation*}
A \cos (2 \pi \omega t+\phi) \tag{1.6}
\end{equation*}
$$

where $A$ is the amplitude, $\omega$ is the frequency of oscillation, and $\phi$ is a phase shift. In (1.5), $A=2, \omega=1 / 50$ (one cycle every 50 time points), and $\phi=.6 \pi$.

An additive noise term was taken to be white noise with $\sigma_{w}=1$ (middle panel) and $\sigma_{w}=5$ (bottom panel), drawn from a normal distribution. Adding the two together obscures the signal, as shown in the lower panels of Figure 1.10. Of course, the degree to which the signal is obscured depends on the amplitude of the signal and the size of $\sigma_{w}$. The ratio of the amplitude of the signal to $\sigma_{w}$ (or some function of the ratio) is sometimes called the signal-to-noise ratio (SNR); the larger the SNR, the easier it is to detect the signal. Note that the signal is easily discernible in the middle panel of Figure 1.10, whereas the signal is obscured in the bottom panel. Typically, we will not observe the signal but the signal obscured by noise.

To reproduce Figure 1.10 in R, use the following commands:

```
cs = 2*cos(2*pi*1:500/50 + .6*pi)
w = rnorm(500,0,1)
par(mfrow=c (3,1), mar=c (3,2,2,1), cex.main=1.5)
```



Fig. 1.10. Cosine wave with period 50 points (top panel) compared with the cosine wave contaminated with additive white Gaussian noise, $\sigma_{w}=1$ (middle panel) and $\sigma_{w}=5$ (bottom panel); see (1.5).

```
plot.ts(cs, main=expression(2*cos(2*pi*t/50+.6*pi)))
plot.ts(cs+w, main=expression(2*cos(2*pi*t/50+.6*pi) + N(0,1)))
plot.ts(cs+5*w, main=expression(2*cos(2*pi*t/50+.6*pi) + N(0,25)))
```


### 1.4 Measures of Dependence

We now discuss various measures that describe the general behavior of a process as it evolves over time. A rather simple descriptive measure is the mean function.

Definition 1.1 The mean function is defined as

$$
\begin{equation*}
\mu_{x t}=E\left(x_{t}\right) \tag{1.7}
\end{equation*}
$$

provided it exists, where E denotes the usual expected value operator. ${ }^{1}$ When no confusion exists about which time series we are referring to, we will drop a subscript and write $\mu_{x t}$ as $\mu_{t}$.

[^1]
## Example 1.11 Mean Function of a Moving Average Series

If $w_{t}$ denotes a white noise series, then $\mu_{w t}=E\left(w_{t}\right)=0$ for all $t$. The top series in Figure 1.7 reflects this, as the series clearly fluctuates around a mean value of zero. Smoothing the series as in Example 1.7 does not change the mean because we can write

$$
\mu_{v t}=E\left(v_{t}\right)=\frac{1}{3}\left[E\left(w_{t-1}\right)+E\left(w_{t}\right)+E\left(w_{t+1}\right)\right]=0 .
$$

## Example 1.12 Mean Function of a Random Walk with Drift

Consider the random walk with drift model given in (1.4),

$$
x_{t}=\delta t+\sum_{j=1}^{t} w_{j}, \quad t=1,2, \ldots
$$

Because $E\left(w_{t}\right)=0$ for all $t$, and $\delta$ is a constant, we have

$$
\mu_{x t}=E\left(x_{t}\right)=\delta t+\sum_{j=1}^{t} E\left(w_{j}\right)=\delta t
$$

which is a straight line with slope $\delta$. A realization of a random walk with drift can be compared to its mean function in Figure 1.9.

## Example 1.13 Mean Function of Signal Plus Noise

A great many practical applications depend on assuming the observed data have been generated by a fixed signal waveform superimposed on a zero-mean noise process, leading to an additive signal model of the form (1.5). It is clear, because the signal in (1.5) is a fixed function of time, we will have

$$
\begin{aligned}
\mu_{x t}=E\left(x_{t}\right) & =E\left[2 \cos \left(2 \pi \frac{t+15}{50}\right)+w_{t}\right] \\
& =2 \cos \left(2 \pi \frac{t+15}{50}\right)+E\left(w_{t}\right) \\
& =2 \cos \left(2 \pi \frac{t+15}{50}\right),
\end{aligned}
$$

and the mean function is just the cosine wave.
The mean function describes only the marginal behavior of a time series. The lack of independence between two adjacent values $x_{s}$ and $x_{t}$ can be assessed numerically, as in classical statistics, using the notions of covariance and correlation. Assuming the variance of $x_{t}$ is finite, we have the following definition.

## Definition 1.2 The autocovariance function is defined as the second moment product

$$
\begin{equation*}
\gamma_{x}(s, t)=\operatorname{cov}\left(x_{s}, x_{t}\right)=E\left[\left(x_{s}-\mu_{s}\right)\left(x_{t}-\mu_{t}\right)\right], \tag{1.8}
\end{equation*}
$$

for all $s$ and $t$. When no possible confusion exists about which time series we are referring to, we will drop the subscript and write $\gamma_{x}(s, t)$ as $\gamma(s, t)$.

Note that $\gamma_{x}(s, t)=\gamma_{x}(t, s)$ for all time points $s$ and $t$. The autocovariance measures the linear dependence between two points on the same series observed at different times. Recall from classical statistics that if $\gamma_{x}(s, t)=0, x_{s}$ and $x_{t}$ are not linearly related, but there still may be some dependence structure between them. If, however, $x_{s}$ and $x_{t}$ are bivariate normal, $\gamma_{x}(s, t)=0$ ensures their independence. It is clear that, for $s=t$, the autocovariance reduces to the (assumed finite) variance, because

$$
\begin{equation*}
\gamma_{x}(t, t)=E\left[\left(x_{t}-\mu_{t}\right)^{2}\right]=\operatorname{var}\left(x_{t}\right) . \tag{1.9}
\end{equation*}
$$

## Example 1.14 Autocovariance of White Noise

The white noise series $w_{t}$ has $E\left(w_{t}\right)=0$ and

$$
\gamma_{w}(s, t)=\operatorname{cov}\left(w_{s}, w_{t}\right)= \begin{cases}\sigma_{w}^{2} & s=t  \tag{1.10}\\ 0 & s \neq t\end{cases}
$$

A realization of white noise with $\sigma_{w}^{2}=1$ is shown in the top panel of Figure 1.7.
We often have to calculate the autocovariance between filtered series. A useful result is given in the following proposition.

Property 1.1 If the random variables

$$
U=\sum_{j=1}^{m} a_{j} X_{j} \quad \text { and } \quad V=\sum_{k=1}^{r} b_{k} Y_{k}
$$

are linear filters of (finite variance) random variables $\left\{X_{j}\right\}$ and $\left\{Y_{k}\right\}$, respectively, then

$$
\begin{equation*}
\operatorname{cov}(U, V)=\sum_{j=1}^{m} \sum_{k=1}^{r} a_{j} b_{k} \operatorname{cov}\left(X_{j}, Y_{k}\right) \tag{1.11}
\end{equation*}
$$

Furthermore, $\operatorname{var}(U)=\operatorname{cov}(U, U)$.
An easy way to remember (1.11) is to treat it like multiplication:

$$
\left(a_{1} X_{1}+a_{2} X_{2}\right)\left(b_{1} Y_{1}\right)=a_{1} b_{1} X_{1} Y_{1}+a_{2} b_{1} X_{2} Y_{1} .
$$

## Example 1.15 Autocovariance of a Moving Average

Consider applying a three-point moving average to the white noise series $w_{t}$ of the previous example as in Example 1.7. In this case,

$$
\gamma_{v}(s, t)=\operatorname{cov}\left(v_{s}, v_{t}\right)=\operatorname{cov}\left\{\frac{1}{3}\left(w_{s-1}+w_{s}+w_{s+1}\right), \frac{1}{3}\left(w_{t-1}+w_{t}+w_{t+1}\right)\right\} .
$$

When $s=t$ we have

$$
\begin{aligned}
\gamma_{v}(t, t) & =\frac{1}{9} \operatorname{cov}\left\{\left(w_{t-1}+w_{t}+w_{t+1}\right),\left(w_{t-1}+w_{t}+w_{t+1}\right)\right\} \\
& =\frac{1}{9}\left[\operatorname{cov}\left(w_{t-1}, w_{t-1}\right)+\operatorname{cov}\left(w_{t}, w_{t}\right)+\operatorname{cov}\left(w_{t+1}, w_{t+1}\right)\right] \\
& =\frac{3}{9} \sigma_{w}^{2} .
\end{aligned}
$$

When $s=t+1$,

$$
\begin{aligned}
\gamma_{v}(t+1, t) & =\frac{1}{9} \operatorname{cov}\left\{\left(w_{t}+w_{t+1}+w_{t+2}\right),\left(w_{t-1}+w_{t}+w_{t+1}\right)\right\} \\
& =\frac{1}{9}\left[\operatorname{cov}\left(w_{t}, w_{t}\right)+\operatorname{cov}\left(w_{t+1}, w_{t+1}\right)\right] \\
& =\frac{2}{9} \sigma_{w}^{2},
\end{aligned}
$$

using (1.10). Similar computations give $\gamma_{v}(t-1, t)=2 \sigma_{w}^{2} / 9, \gamma_{v}(t+2, t)=\gamma_{v}(t-$ $2, t)=\sigma_{w}^{2} / 9$, and 0 when $|t-s|>2$. We summarize the values for all $s$ and $t$ as

$$
\gamma_{v}(s, t)= \begin{cases}\frac{3}{9} \sigma_{w}^{2} & s=t  \tag{1.12}\\ \frac{2}{9} \sigma_{w}^{2} & |s-t|=1 \\ \frac{1}{9} \sigma_{w}^{2} & |s-t|=2 \\ 0 & |s-t|>2\end{cases}
$$

Example 1.15 shows clearly that the smoothing operation introduces a covariance function that decreases as the separation between the two time points increases and disappears completely when the time points are separated by three or more time points. This particular autocovariance is interesting because it only depends on the time separation or lag and not on the absolute location of the points along the series. We shall see later that this dependence suggests a mathematical model for the concept of weak stationarity.

## Example 1.16 Autocovariance of a Random Walk

For the random walk model, $x_{t}=\sum_{j=1}^{t} w_{j}$, we have

$$
\gamma_{x}(s, t)=\operatorname{cov}\left(x_{s}, x_{t}\right)=\operatorname{cov}\left(\sum_{j=1}^{s} w_{j}, \sum_{k=1}^{t} w_{k}\right)=\min \{s, t\} \sigma_{w}^{2}
$$

because the $w_{t}$ are uncorrelated random variables. Note that, as opposed to the previous examples, the autocovariance function of a random walk depends on the particular time values $s$ and $t$, and not on the time separation or lag. Also, notice that the variance of the random walk, $\operatorname{var}\left(x_{t}\right)=\gamma_{x}(t, t)=t \sigma_{w}^{2}$, increases without bound as time $t$ increases. The effect of this variance increase can be seen in Figure 1.9 where the processes start to move away from their mean functions $\delta t$ (note that $\delta=0$ and .2 in that example).

As in classical statistics, it is more convenient to deal with a measure of association between -1 and 1 , and this leads to the following definition.

Definition 1.3 The autocorrelation function (ACF) is defined as

$$
\begin{equation*}
\rho(s, t)=\frac{\gamma(s, t)}{\sqrt{\gamma(s, s) \gamma(t, t)}} . \tag{1.13}
\end{equation*}
$$

The ACF measures the linear predictability of the series at time $t$, say $x_{t}$, using only the value $x_{s}$. We can show easily that $-1 \leq \rho(s, t) \leq 1$ using the Cauchy-Schwarz inequality. ${ }^{2}$ If we can predict $x_{t}$ perfectly from $x_{s}$ through a linear relationship, $x_{t}=\beta_{0}+\beta_{1} x_{s}$, then the correlation will be +1 when $\beta_{1}>0$, and -1 when $\beta_{1}<0$. Hence, we have a rough measure of the ability to forecast the series at time $t$ from the value at time $s$.

Often, we would like to measure the predictability of another series $y_{t}$ from the series $x_{s}$. Assuming both series have finite variances, we have the following definition.

Definition 1.4 The cross-covariance function between two series, $x_{t}$ and $y_{t}$, is

$$
\begin{equation*}
\gamma_{x y}(s, t)=\operatorname{cov}\left(x_{s}, y_{t}\right)=E\left[\left(x_{s}-\mu_{x s}\right)\left(y_{t}-\mu_{y t}\right)\right] . \tag{1.14}
\end{equation*}
$$

The cross-covariance function can be scaled to live in $[-1,1]$ :
Definition 1.5 The cross-correlation function (CCF) is given by

$$
\begin{equation*}
\rho_{x y}(s, t)=\frac{\gamma_{x y}(s, t)}{\sqrt{\gamma_{x}(s, s) \gamma_{y}(t, t)}} . \tag{1.15}
\end{equation*}
$$

### 1.5 Stationary Time Series

The preceding definitions of the mean and autocovariance functions are completely general. Although we have not made any special assumptions about the behavior of the time series, many of the preceding examples have hinted that a sort of regularity may exist over time in the behavior of a time series.

Definition 1.6 A strictly stationary time series is one for which the probabilistic behavior of every collection of values

$$
\left\{x_{t_{1}}, x_{t_{2}}, \ldots, x_{t_{k}}\right\}
$$

is identical to that of the time shifted set

$$
\left\{x_{t_{1}+h}, x_{t_{2}+h}, \ldots, x_{t_{k}+h}\right\}
$$

for all $k=1,2, \ldots$, all time points $t_{1}, t_{2}, \ldots, t_{k}$, and all time shifts $h=0, \pm 1, \pm 2, \ldots$.
It is difficult to assess strict stationarity from data. Rather than imposing conditions on all possible distributions of a time series, we will use a milder version that imposes conditions only on the first two moments of the series.

[^2]Definition 1.7 A weakly stationary time series is a finite variance process where
(i) the mean value function, $\mu_{t}$, defined in (1.7) is constant and does not depend on time $t$, and
(ii) the autocovariance function, $\gamma(s, t)$, defined in (1.8) depends on $s$ and $t$ only through their difference $|s-t|$.

Henceforth, we will use the term stationary to mean weakly stationary; if a process is stationary in the strict sense, we will use the term strictly stationary.

Stationarity requires regularity in the mean and autocorrelation functions so that these quantities (at least) may be estimated by averaging. It should be clear that a strictly stationary, finite variance, time series is also stationary. The converse is not true in general. One important case where stationarity implies strict stationarity is if the time series is Gaussian [meaning all finite collections of the series are Gaussian].

## Example 1.17 A Random Walk is Not Stationary

A random walk is not stationary because its autocovariance function, $\gamma(s, t)=$ $\min \{s, t\} \sigma_{w}^{2}$, depends on time; see Example 1.16 and Problem 1.6. Also, the random walk with drift violates both conditions of Definition 1.7 because, as shown in Example 1.12, the mean function, $\mu_{x t}=\delta t$, is also a function of time $t$.

Because the mean function, $E\left(x_{t}\right)=\mu_{t}$, of a stationary time series is independent of time $t$, we will write

$$
\begin{equation*}
\mu_{t}=\mu \tag{1.16}
\end{equation*}
$$

Also, because the autocovariance function, $\gamma(s, t)$, of a stationary time series, $x_{t}$, depends on $s$ and $t$ only through their difference $|s-t|$, we may simplify the notation. Let $s=t+h$, where $h$ represents the time shift or lag. Then

$$
\gamma(t+h, t)=\operatorname{cov}\left(x_{t+h}, x_{t}\right)=\operatorname{cov}\left(x_{h}, x_{0}\right)=\gamma(h, 0)
$$

because the time difference between times $t+h$ and $t$ is the same as the time difference between times $h$ and 0 . Thus, the autocovariance function of a stationary time series does not depend on the time argument $t$. Henceforth, for convenience, we will drop the second argument of $\gamma(h, 0)$.

Definition 1.8 The autocovariance function of a stationary time series will be written as

$$
\begin{equation*}
\gamma(h)=\operatorname{cov}\left(x_{t+h}, x_{t}\right)=E\left[\left(x_{t+h}-\mu\right)\left(x_{t}-\mu\right)\right] . \tag{1.17}
\end{equation*}
$$

Definition 1.9 The autocorrelation function (ACF) of a stationary time series will be written using (1.13) as

$$
\begin{equation*}
\rho(h)=\frac{\gamma(h)}{\gamma(0)} . \tag{1.18}
\end{equation*}
$$

The Cauchy-Schwarz inequality shows again that $-1 \leq \rho(h) \leq 1$ for all $h$, enabling one to assess the relative importance of a given autocorrelation value by comparing with the extreme values -1 and 1 .


Fig. 1.11. Autocovariance function of a three-point moving average.

## Example 1.18 Stationarity of White Noise

The mean and autocovariance functions of the white noise series discussed in Example 1.6 and Example 1.14 are easily evaluated as $\mu_{w t}=0$ and

$$
\gamma_{w}(h)=\operatorname{cov}\left(w_{t+h}, w_{t}\right)= \begin{cases}\sigma_{w}^{2} & h=0 \\ 0 & h \neq 0\end{cases}
$$

Thus, white noise satisfies the conditions of Definition 1.7 and is weakly stationary or stationary.

## Example 1.19 Stationarity of a Moving Average

The three-point moving average process of Example 1.7 is stationary because, from Example 1.11 and Example 1.15, the mean and autocovariance functions $\mu_{v t}=0$, and

$$
\gamma_{v}(h)= \begin{cases}\frac{3}{9} \sigma_{w}^{2} & h=0 \\ \frac{2}{9} \sigma_{w}^{2} & h= \pm 1 \\ \frac{1}{9} \sigma_{w}^{2} & h= \pm 2 \\ 0 & |h|>2\end{cases}
$$

are independent of time $t$, satisfying the conditions of Definition 1.7.
Note that the ACF is given by

$$
\rho_{v}(h)=\left\{\begin{array}{ll}
1 & h=0, \\
\frac{2}{3} & h= \pm 1, \\
\frac{1}{3} & h= \pm 2, \\
0 & |h|>2
\end{array} .\right.
$$

Figure 1.11 shows a plot of the autocorrelation as a function of lag $h$. Note that the autocorrelation function is symmetric about lag zero and in this example, decays as a function of lag.

## Example 1.20 Trend Stationarity

If $x_{t}=\alpha+\beta t+w_{t}$, then the mean function is $\mu_{x, t}=E\left(x_{t}\right)=\alpha+\beta t$, which is not independent of time. Therefore, the process is not stationary. The autocovariance function, however, is independent of time, because $\gamma_{x}(h)=\operatorname{cov}\left(x_{t+h}, x_{t}\right)=$ $E\left[\left(x_{t+h}-\mu_{x, t+h}\right)\left(x_{t}-\mu_{x, t}\right)\right]=E\left(w_{t+h} w_{t}\right)=\gamma_{w}(h)$, which is given in Example 1.18. Thus, the model may be considered as having stationary behavior around a linear trend; this behavior is sometimes called trend stationarity.

The autocovariance function of a stationary process has several useful properties. First, the value at $h=0$, namely

$$
\begin{equation*}
\gamma(0)=E\left[\left(x_{t}-\mu\right)^{2}\right]=\operatorname{var}\left(x_{t}\right) . \tag{1.19}
\end{equation*}
$$

Also, the Cauchy-Schwarz inequality implies $|\gamma(h)| \leq \gamma(0)$. Another useful property is that the autocovariance function of a stationary series is symmetric around the origin,

$$
\begin{equation*}
\gamma(h)=\gamma(-h) \tag{1.20}
\end{equation*}
$$

for all $h$. This property follows because

$$
\begin{aligned}
\gamma(h) & =\gamma((t+h)-t)=E\left[\left(x_{t+h}-\mu\right)\left(x_{t}-\mu\right)\right] \\
& =E\left[\left(x_{t}-\mu\right)\left(x_{t+h}-\mu\right)\right]=\gamma(t-(t+h))=\gamma(-h),
\end{aligned}
$$

which shows how to use the notation as well as proving the result.
When several series are available, a notion of stationarity still applies with additional conditions.

Definition 1.10 Two time series, say, $x_{t}$ and $y_{t}$, are said to be jointly stationary if they are each stationary, and the cross-covariance function

$$
\begin{equation*}
\gamma_{x y}(h)=\operatorname{cov}\left(x_{t+h}, y_{t}\right)=E\left[\left(x_{t+h}-\mu_{x}\right)\left(y_{t}-\mu_{y}\right)\right] \tag{1.21}
\end{equation*}
$$

is a function only of lag $h$.
Definition 1.11 The cross-correlation function (CCF) of jointly stationary time series $x_{t}$ and $y_{t}$ is defined as

$$
\begin{equation*}
\rho_{x y}(h)=\frac{\gamma_{x y}(h)}{\sqrt{\gamma_{x}(0) \gamma_{y}(0)}} . \tag{1.22}
\end{equation*}
$$

Again, we have the result $-1 \leq \rho_{x y}(h) \leq 1$ which enables comparison with the extreme values -1 and 1 when looking at the relation between $x_{t+h}$ and $y_{t}$. The cross-correlation function is not generally symmetric about zero [i.e., typically $\left.\rho_{x y}(h) \neq \rho_{x y}(-h)\right]$; however, it is the case that

$$
\begin{equation*}
\rho_{x y}(h)=\rho_{y x}(-h), \tag{1.23}
\end{equation*}
$$

which can be shown by manipulations similar to those used to show (1.20).


Fig. 1.12. Demonstration of the results of Example 1.22 when $\ell=5$.

## Example 1.21 Joint Stationarity

Consider the two series, $x_{t}$ and $y_{t}$, formed from the sum and difference of two successive values of a white noise process, say,

$$
x_{t}=w_{t}+w_{t-1} \quad \text { and } \quad y_{t}=w_{t}-w_{t-1},
$$

where $w_{t}$ are independent random variables with zero means and variance $\sigma_{w}^{2}$. It is easy to show that $\gamma_{x}(0)=\gamma_{y}(0)=2 \sigma_{w}^{2}$ and $\gamma_{x}(1)=\gamma_{x}(-1)=\sigma_{w}^{2}, \gamma_{y}(1)=$ $\gamma_{y}(-1)=-\sigma_{w}^{2}$. Also,

$$
\gamma_{x y}(1)=\operatorname{cov}\left(x_{t+1}, y_{t}\right)=\operatorname{cov}\left(w_{t+1}+w_{t}, w_{t}-w_{t-1}\right)=\sigma_{w}^{2}
$$

because only one term is nonzero (recall Property 1.1 on page 12). Similarly, $\gamma_{x y}(0)=0, \gamma_{x y}(-1)=-\sigma_{w}^{2}$. We obtain, using (1.22),

$$
\rho_{x y}(h)=\left\{\begin{array}{cl}
0 & h=0 \\
1 / 2 & h=1 \\
-1 / 2 & h=-1 \\
0 & |h| \geq 2
\end{array}\right.
$$

Clearly, the autocovariance and cross-covariance functions depend only on the lag separation, $h$, so the series are jointly stationary.

## Example 1.22 Prediction Using Cross-Correlation

Consider the problem of determining possible leading or lagging relations between two series $x_{t}$ and $y_{t}$. If the model

$$
y_{t}=A x_{t-\ell}+w_{t}
$$

holds, the series $x_{t}$ is said to lead $y_{t}$ for $\ell>0$ and is said to lag $y_{t}$ for $\ell<0$. Hence, the analysis of leading and lagging relations might be important in predicting the value of $y_{t}$ from $x_{t}$. Assuming that the noise $w_{t}$ is uncorrelated with the $x_{t}$ series, the cross-covariance function can be computed as

$$
\begin{aligned}
\gamma_{x y}(h) & =\operatorname{cov}\left(x_{t+h}, y_{t}\right)=\operatorname{cov}\left(x_{t+h}, A x_{t-\ell}+w_{t}\right) \\
& =\operatorname{cov}\left(x_{t+h}, A x_{t-\ell}\right)=A \gamma_{x}(h+\ell) .
\end{aligned}
$$

Since the largest value of $\gamma_{x}(h+\ell)$ is $\gamma_{x}(0)$, i.e., when $h=-\ell$, the cross-covariance function will look like the autocovariance of the input series $x_{t}$, and it will have a peak on the negative side if $x_{t}$ leads $y_{t}$ and a peak on the positive side if $x_{t}$ lags $y_{t}$. Below is the R code of an example with $\ell=5$ and $\hat{\gamma}_{x y}(h)$ is shown in Figure 1.12. set. seed $(90210) ; \quad x=\operatorname{rnorm}(100) ; \quad y=\operatorname{lag}(x,-5)+\operatorname{rnorm}(100)$ $\operatorname{ccf}(x, y, y l a b=' C \operatorname{lovF',~type='covariance')~}$

### 1.6 Estimation of Correlation

For data analysis, only the sample values, $x_{1}, x_{2}, \ldots, x_{n}$, are available for estimating the mean, autocovariance, and autocorrelation functions. In this case, the assumption of stationarity becomes critical and allows the use of averaging to estimate the population means and covariance functions.

Accordingly, if a time series is stationary, the mean function (1.16) $\mu_{t}=\mu$ is constant so that we can estimate it by the sample mean,

$$
\begin{equation*}
\bar{x}=\frac{1}{n} \sum_{t=1}^{n} x_{t} . \tag{1.24}
\end{equation*}
$$

The standard error of the estimate is the square root of $\operatorname{var}(\bar{x})$, which can be computed using first principles (recall Property 1.1), and is given by

$$
\begin{equation*}
\operatorname{var}(\bar{x})=\frac{1}{n^{2}} \operatorname{cov}\left(\sum_{t=1}^{n} x_{t}, \sum_{s=1}^{n} x_{s}\right)=\frac{1}{n^{2}} \sum_{t=1}^{n} \sum_{s=1}^{n} \operatorname{cov}\left(x_{t}, x_{s}\right)=\frac{1}{n} \sum_{h=-n}^{n}\left(1-\frac{|h|}{n}\right) \gamma_{x}(h) ; \tag{1.25}
\end{equation*}
$$

see Figure 1.13.


Fig. 1.13. Summing over $s$ and $t$ in (1.25) with $n=5$.

If the process is white noise, (1.25) reduces to the familiar $\sigma_{x}^{2} / n$ recalling that $\gamma_{x}(0)=\sigma_{x}^{2}$. Note that, in the case of dependence, the standard error of $\bar{x}$ may be
smaller or larger than the white noise case depending on the nature of the correlation structure (see Problem 1.13).

The theoretical autocovariance function, (1.17), is estimated by the sample autocovariance function defined as follows.

Definition 1.12 The sample autocovariance function is defined as

$$
\begin{equation*}
\widehat{\gamma}(h)=n^{-1} \sum_{t=1}^{n-h}\left(x_{t+h}-\bar{x}\right)\left(x_{t}-\bar{x}\right), \tag{1.26}
\end{equation*}
$$

with $\widehat{\gamma}(-h)=\widehat{\gamma}(h)$ for $h=0,1, \ldots, n-1$.
The sum in (1.26) runs over a restricted range because $x_{t+h}$ is not available for $t+h>n$. The estimator in (1.26) is preferred to the one that would be obtained by dividing by $n-h$ because (1.26) is a non-negative definite function.

The autocovariance function, $\gamma(h)$, of a stationary process is non-negative definite ensuring that variances of linear combinations of the variates $x_{t}$ will never be negative. That is, for any $n \geq 1$, and constants $a_{1}, \ldots, a_{n}$,

$$
0 \leq \operatorname{var}\left(a_{1} x_{1}+\cdots+a_{n} x_{n}\right)=\sum_{j=1}^{n} \sum_{k=1}^{n} a_{j} a_{k} \gamma(j-k)
$$

using Property 1.1. And, because a variance is never negative, the estimate of that variance

$$
\widehat{\operatorname{var}}\left(a_{1} x_{1}+\cdots+a_{n} x_{n}\right)=\sum_{j=1}^{n} \sum_{k=1}^{n} a_{j} a_{k} \widehat{\gamma}(j-k)
$$

should also be non-negative. The estimator in (1.26) guarantees this result, but no such guarantee exists if we divide by $n-h$. Note that neither dividing by $n$ nor $n-h$ in (1.26) yields an unbiased estimator of $\gamma(h)$.

Definition 1.13 The sample autocorrelation function is defined, analogously to (1.18), as

$$
\begin{equation*}
\widehat{\rho}(h)=\frac{\widehat{\gamma}(h)}{\widehat{\gamma}(0)} . \tag{1.27}
\end{equation*}
$$

## Example 1.23 Sample ACF and Scatterplots

Estimating autocorrelation is similar to estimating of correlation in the usual setup where we have pairs of observations, say $\left(x_{i}, y_{i}\right)$, for $i=1, \ldots, n$. For example, if we have time series data $x_{t}$ for $t=1, \ldots, n$, then the pairs of observations for estimating $\rho(h)$ are the $n-h$ pairs given by $\left\{\left(x_{t}, x_{t+h}\right) ; t=1, \ldots, n-h\right\}$. Figure 1.14 shows an example using the SOI series where $\widehat{\rho}(1)=.604$ and $\widehat{\rho}(6)=-.187$. The following code was used for Figure 1.14.



Fig. 1.14. Display for Example 1.23. For the SOI series, we have a scatterplot of pairs of values one month apart (left) and six months apart (right). The estimated correlation is displayed in the box.

```
(r = round(acf(soi, 6, plot=FALSE)$acf[-1], 3)) # first 6 sample acf values
    [1] 0.604 0.374 0.214 0.050 -0.107 -0.187
par(mfrow=c(1,2), mar=c (3,3,1,1), mgp=c(1.6,.6,0))
plot(lag(soi,-1), soi)
legend('topleft', legend=r[1])
plot(lag(soi,-6), soi)
legend('topleft', legend=r[6])
```

The sample autocorrelation function has a sampling distribution that allows us to assess whether the data comes from a completely random or white series or whether correlations are statistically significant at some lags.

## Property 1.2 Large-Sample Distribution of the ACF

If $x_{t}$ is white noise, then for $n$ large and under mild conditions, the sample ACF, $\widehat{\rho}_{x}(h)$, for $h=1,2, \ldots, H$, where $H$ is fixed but arbitrary, is approximately normal with zero mean and standard deviation given by of $\frac{1}{\sqrt{n}}$.

Based on Property 1.2, we obtain a rough method for assessing whether a series is white noise by determining how many values of $\widehat{\rho}(h)$ are outside the interval $\pm 2 / \sqrt{n}$ (two standard errors); for white noise, approximately $95 \%$ of the sample ACFs should be within these limits. ${ }^{3}$ The bounds do not hold in general and can be ignored if the interest is other than assessing whiteness. The applications of this property develop because many statistical modeling procedures depend on reducing a time series to a white noise series using various kinds of transformations. Afterwards the plotted ACF of the residuals behave as stated.

## Example 1.24 A Simulated Time Series

To compare the sample ACF for various sample sizes to the theoretical ACF, consider a contrived set of data generated by tossing a fair coin, letting $x_{t}=1$ when a head is obtained and $x_{t}=-1$ when a tail is obtained. Then, construct $y_{t}$ as

[^3]

Fig. 1.15. Realization of (1.28), $n=10$.

$$
\begin{equation*}
y_{t}=5+x_{t}-.7 x_{t-1} \tag{1.28}
\end{equation*}
$$

To simulate data, we consider two cases, one with a small sample size ( $n=10$; see Figure 1.15) and another with a moderate sample size $(n=100)$.

```
set.seed(101010)
x1 = 2*rbinom(11, 1, .5) - 1 # simulated sequence of coin tosses
x2 = 2*rbinom(101, 1, .5) - 1
y1 = 5 + filter(x1, sides=1, filter=c(1,-.7))[-1]
y2 = 5 + filter(x2, sides=1, filter=c(1,-.7))[-1]
plot.ts(y1, type='s'); plot.ts(y2, type='s') # only one shown
acf(y1, lag.max=4, plot=FALSE) # 1/\sqrt{}{10}=.32
    Autocorrelations of series 'y1', by lag
            0
    1.000-0.688 0.425-0.306-0.007
acf(y2, lag.max=4, plot=FALSE) # 1/\sqrt{}{100}=.1
    Autocorrelations of series 'y2', by lag
            0
    1.000-0.480-0.002 -0.004 0.000
# Note that the sample ACF at lag zero is always 1 (Why?).
```

The theoretical ACF can be obtained from the model (1.28) using first principles so that

$$
\rho_{y}(1)=\frac{-.7}{1+.7^{2}}=-.47
$$

and $\rho_{y}(h)=0$ for $|h|>1$ (Problem 1.18). It is interesting to compare the theoretical ACF with sample ACFs for the realization where $n=10$ and the other realization where $n=100$; note the increased variability in the smaller size sample.

Definition 1.14 The estimators for the cross-covariance function, $\gamma_{x y}(h)$, as given in (1.21) and the cross-correlation, $\rho_{x y}(h)$, in (1.22) are given, respectively, by the sample cross-covariance function

$$
\begin{equation*}
\widehat{\gamma}_{x y}(h)=n^{-1} \sum_{t=1}^{n-h}\left(x_{t+h}-\bar{x}\right)\left(y_{t}-\bar{y}\right), \tag{1.29}
\end{equation*}
$$

where $\widehat{\gamma}_{x y}(-h)=\widehat{\gamma}_{y x}(h)$ determines the function for negative lags, and the sample cross-correlation function


Fig. 1.16. Sample ACFs of the SOI series (top) and of the Recruitment series (middle), and the sample CCF of the two series (bottom); negative lags indicate SOI leads Recruitment. The lag axes are in terms of seasons ( 12 months).

$$
\begin{equation*}
\widehat{\rho}_{x y}(h)=\frac{\widehat{\gamma}_{x y}(h)}{\sqrt{\widehat{\gamma}_{x}(0) \widehat{\gamma}_{y}(0)}} . \tag{1.30}
\end{equation*}
$$

The sample cross-correlation function can be examined graphically as a function of lag $h$ to search for leading or lagging relations in the data using the property mentioned in Example 1.22 for the theoretical cross-covariance function. Because $-1 \leq \widehat{\rho}_{x y}(h) \leq 1$, the practical importance of peaks can be assessed by comparing their magnitudes with their theoretical maximum values. Furthermore, for $x_{t}$ and $y_{t}$ independent processes, then under mild conditions, we have the following property.

## Property 1.3 Large-Sample Distribution of Cross-Correlation

The large sample distribution of $\widehat{\rho}_{x y}(h)$ is normal with mean zero and standard deviation $\frac{1}{\sqrt{n}}$ if at least one of the processes is independent white noise.

## Example 1.25 SOI and Recruitment Correlation Analysis

The autocorrelation and cross-correlation functions are also useful for analyzing the joint behavior of two stationary series whose behavior may be related in some unspecified way. In Example 1.4 (see Figure 1.5), we have considered simultaneous monthly readings of the SOI and the number of new fish (Recruitment) computed
from a model. Figure 1.16 shows the autocorrelation and cross-correlation functions (ACFs and CCF) for these two series.

Both of the ACFs exhibit periodicities corresponding to the correlation between values separated by 12 units. Observations 12 months or one year apart are strongly positively correlated, as are observations at multiples such as $24,36,48, \ldots$ Observations separated by six months are negatively correlated, showing that positive excursions tend to be associated with negative excursions six months removed. This appearance is rather characteristic of the pattern that would be produced by a sinusoidal component with a period of 12 months; see Example 1.26. The crosscorrelation function peaks at $h=-6$, showing that the SOI measured at time $t-6$ months is associated with the Recruitment series at time $t$. We could say the SOI leads the Recruitment series by six months. The sign of the CCF at $h=-6$ is negative, leading to the conclusion that the two series move in different directions; that is, increases in SOI lead to decreases in Recruitment and vice versa. Again, note the periodicity of 12 months in the CCF.

The flat lines shown on the plots indicate $\pm 2 / \sqrt{453}$, so that upper values would be exceeded about $2.5 \%$ of the time if the noise were white as specified in Property 1.2 and Property 1.3. Of course, neither series is noise, so we can ignore these lines. To reproduce Figure 1.16 in R, use the following commands:
$\operatorname{par}(m f r o w=c(3,1))$
acf(soi, 48, main="Southern Oscillation Index")
acf(rec, 48, main="Recruitment")
ccf(soi, rec, 48, main="SOI vs Recruitment", ylab="CCF")

## Example 1.26 Prewhitening and Cross Correlation Analysis

Although we do not have all the tools necessary yet, it is worthwhile to discuss the idea of prewhitening a series prior to a cross-correlation analysis. The basic idea is simple; in order to use Property 1.3, at least one of the series must be white noise. If this is not the case, there is no simple way to tell if a cross-correlation estimate is significantly different from zero. Hence, in Example 1.25, we were only guessing at the linear dependence relationship between SOI and Recruitment.

For example, in Figure 1.17 we generated two series, $x_{t}$ and $y_{t}$, for $t=1, \ldots, 120$ independently as

$$
x_{t}=2 \cos \left(2 \pi t \frac{1}{12}\right)+w_{t 1} \quad \text { and } \quad y_{t}=2 \cos \left(2 \pi[t+5] \frac{1}{12}\right)+w_{t 2}
$$

where $\left\{w_{t 1}, w_{t 2} ; t=1, \ldots, 120\right\}$ are all independent standard normals. The series are made to resemble SOI and Recruitment. The generated data are shown in the top row of the figure. The middle row of Figure 1.17 show the sample ACF of each series, each of which exhibits the cyclic nature of each series. The bottom row (left) of Figure 1.17 shows the sample CCF between $x_{t}$ and $y_{t}$, which appears to show cross-correlation even though the series are independent. The bottom row (right) also displays the sample CCF between $x_{t}$ and the prewhitened $y_{t}$, which shows that the two sequences are uncorrelated. By prewhtiening $y_{t}$, we mean that the signal has been removed from the data by running a regression of $y_{t}$ on $\cos (2 \pi t)$ and $\sin (2 \pi t)$ [see Example 2.9] and then putting $\tilde{y}_{t}=y_{t}-\hat{y}_{t}$, where $\hat{y}_{t}$ are the predicted values from the regression.


Fig. 1.17. Display for Example 1.26

The following code will reproduce Figure 1.17.

```
set.seed(1492)
num=120; t=1:num
X = ts(2*cos(2*pi*t/12) + rnorm(num), freq=12)
Y = ts(2*}\operatorname{cos(2*pi*(t+5)/12) + rnorm(num), freq=12)
Yw = resid( lm(Y~ cos(2*pi*t/12) + sin(2*pi*t/12), na.action=NULL) )
par(mfrow=c (3,2), mgp=c(1.6,.6,0), mar=c (3,3,1,1) )
plot(X); plot(Y)
acf(X,48, ylab='ACF(X)'); acf(Y,48, ylab='ACF(Y)')
ccf(X,Y,24, ylab='CCF(X,Y)')
ccf(X,Yw,24, ylab='CCF(X,Yw)', ylim=c(-.6,.6))
```


## Problems

1.1 In 25 words or less, and without using symbols, why is stationarity important?
1.2 (a) Generate $n=100$ observations from the autoregression

$$
x_{t}=-.9 x_{t-2}+w_{t}
$$

with $\sigma_{w}=1$, using the method described in Example 1.8. Next, apply the moving average filter

$$
v_{t}=\left(x_{t}+x_{t-1}+x_{t-2}+x_{t-3}\right) / 4
$$

to $x_{t}$, the data you generated. Now plot $x_{t}$ as a line and superimpose $v_{t}$ as a dashed line. Note: $\mathrm{v}=\operatorname{filter}(\mathrm{x}, \operatorname{rep}(1 / 4,4), \operatorname{sides}=1)$
(b) Repeat (a) but with

$$
x_{t}=2 \cos (2 \pi t / 4)+w_{t}
$$

where $w_{t} \sim \operatorname{iid} \mathrm{~N}(0,1)$.
(c) Repeat (a) but where $x_{t}$ is the log of the Johnson \& Johnson data discussed in Example 1.1.
(d) What is seasonal adjustment (you can do an internet search)?
(e) State your conclusions (in other words, what did you learn from this exercise).
1.3 Show that the autocovariance function can be written as

$$
\gamma(s, t)=E\left[\left(x_{s}-\mu_{s}\right)\left(x_{t}-\mu_{t}\right)\right]=E\left(x_{s} x_{t}\right)-\mu_{s} \mu_{t},
$$

where $E\left[x_{t}\right]=\mu_{t}$.
1.4 Consider the time series

$$
x_{t}=\beta_{1}+\beta_{2} t+w_{t}
$$

where $\beta_{1}$ and $\beta_{2}$ are regression coefficients, and $w_{t}$ is a white noise process with variance $\sigma_{w}^{2}$.
(a) Determine whether $x_{t}$ is stationary.
(b) Show that the process $y_{t}=x_{t}-x_{t-1}$ is stationary.
(c) Show that the mean of the moving average

$$
v_{t}=\frac{1}{3}\left(x_{t-1}+x_{t}+x_{t+1}\right)
$$

is $\beta_{1}+\beta_{2} t$.
1.5 For a moving average process of the form

$$
x_{t}=w_{t-1}+2 w_{t}+w_{t+1}
$$

where $w_{t}$ are independent with zero means and variance $\sigma_{w}^{2}$, determine the autocovariance and autocorrelation functions as a function of lag $h$ and sketch the ACF as a function of $h$.
1.6 Consider the random walk with drift model

$$
x_{t}=\delta+x_{t-1}+w_{t}
$$

for $t=1,2, \ldots$, with $x_{0}=0$, where $w_{t}$ is white noise with variance $\sigma_{w}^{2}$.
(a) Show that the model can be written as $x_{t}=\delta t+\sum_{k=1}^{t} w_{k}$.
(b) Find the mean function and the autocovariance function of $x_{t}$.
(c) Argue that $x_{t}$ is not stationary.
(d) Show $\rho_{x}(t-1, t)=\sqrt{\frac{t-1}{t}} \rightarrow 1$ as $t \rightarrow \infty$. What is the implication of this result?
(e) Suggest a transformation to make the series stationary, and prove that the transformed series is stationary. (Hint: See Problem 1.4b.)
1.7 Would you treat the global temperature data discussed in Example 1.2 and shown in Figure 1.3 as stationary or non-stationary? Support your answer.
1.8 A time series with a periodic component can be constructed from

$$
x_{t}=U_{1} \sin \left(2 \pi \omega_{0} t\right)+U_{2} \cos \left(2 \pi \omega_{0} t\right)
$$

where $U_{1}$ and $U_{2}$ are independent random variables with zero means and $E\left(U_{1}^{2}\right)=$ $E\left(U_{2}^{2}\right)=\sigma^{2}$. The constant $\omega_{0}$ determines the period or time it takes the process to make one complete cycle. Show that this series is weakly stationary with autocovariance function

$$
\gamma(h)=\sigma^{2} \cos \left(2 \pi \omega_{0} h\right)
$$

1.9 Suppose we would like to predict a single stationary series $x_{t}$ with zero mean and autocorrelation function $\gamma(h)$ at some time in the future, say, $t+m$, for $m>0$.
(a) If we predict using only $x_{t}$ and some scale multiplier $A$, show that the mean-square prediction error

$$
\operatorname{MSE}(A)=E\left[\left(x_{t+m}-A x_{t}\right)^{2}\right]
$$

is minimized by the value

$$
A=\rho(m) .
$$

(b) Show that the minimum mean-square prediction error is

$$
\operatorname{MSE}(A)=\gamma(0)\left[1-\rho^{2}(m)\right]
$$

(c) Show that if $x_{t+m}=A x_{t}$, then $\rho(m)=1$ if $A>0$, and $\rho(m)=-1$ if $A<0$.
1.10 For two jointly stationary series $x_{t}$ and $y_{t}$, verify (1.23).
1.11 Consider the two series

$$
\begin{gathered}
x_{t}=w_{t} \\
y_{t}=w_{t}-\theta w_{t-1}+u_{t}
\end{gathered}
$$

where $w_{t}$ and $u_{t}$ are independent white noise series with variances $\sigma_{w}^{2}$ and $\sigma_{u}^{2}$, respectively, and $\theta$ is an unspecified constant.
(a) Express the ACF, $\rho_{y}(h)$, for $h=0, \pm 1, \pm 2, \ldots$ of the series $y_{t}$ as a function of $\sigma_{w}^{2}, \sigma_{u}^{2}$, and $\theta$.
(b) Determine the CCF, $\rho_{x y}(h)$ relating $x_{t}$ and $y_{t}$.
(c) Show that $x_{t}$ and $y_{t}$ are jointly stationary.
1.12 Let $w_{t}$, for $t=0, \pm 1, \pm 2, \ldots$ be a normal white noise process, and consider the series

$$
x_{t}=w_{t} w_{t-1}
$$

Determine the mean and autocovariance function of $x_{t}$, and state whether it is stationary.
1.13 Suppose $x_{t}=\mu+w_{t}+\theta w_{t-1}$, where $w_{t} \sim w n\left(0, \sigma_{w}^{2}\right)$.
(a) Show that mean function is $E\left(x_{t}\right)=\mu$.
(b) Show that the autocovariance function of $x_{t}$ is given by $\gamma_{x}(0)=\sigma_{w}^{2}\left(1+\theta^{2}\right)$, $\gamma_{x}( \pm 1)=\sigma_{w}^{2} \theta$, and $\gamma_{x}(h)=0$ otherwise.
(c) Show that $x_{t}$ is stationary for all values of $\theta \in \mathbb{R}$.
(d) Use (1.25) to calculate $\operatorname{var}(\bar{x})$ for estimating $\mu$ when (i) $\theta=1$, (ii) $\theta=0$, and (iii) $\theta=-1$
(e) In time series, the sample size $n$ is typically large, so that $\frac{(n-1)}{n} \approx 1$. With this as a consideration, comment on the results of part (d); in particular, how does the accuracy in the estimate of the mean $\mu$ change for the three different cases?
1.14 (a) Simulate a series of $n=500$ Gaussian white noise observations as in Example 1.6 and compute the sample ACF, $\widehat{\rho}(h)$, to lag 20. Compare the sample ACF you obtain to the actual ACF, $\rho(h)$. [Recall Example 1.18.]
(b) Repeat part (a) using only $n=50$. How does changing $n$ affect the results?
1.15 (a) Simulate a series of $n=500$ moving average observations as in Example 1.7 and compute the sample ACF, $\widehat{\rho}(h)$, to lag 20 . Compare the sample ACF you obtain to the actual ACF, $\rho(h)$. [Recall Example 1.19.]
(b) Repeat part (a) using only $n=50$. How does changing $n$ affect the results?
1.16 Simulate 500 observations from the AR model specified in Example 1.8 and then plot the sample ACF to lag 50. What does the sample ACF tell you about the approximate cyclic behavior of the data? Hint: Recall Example 1.25.
1.17 Simulate a series of $n=500$ observations from the signal-plus-noise model presented in Example 1.10 with (a) $\sigma_{w}=0$, (b) $\sigma_{w}=1$ and (c) $\sigma_{w}=5$. Compute the sample ACF to lag 100 of the three series you generated and comment.
1.18 For the time series $y_{t}$ described in Example 1.24, verify the stated result that $\rho_{y}(1)=-.47$ and $\rho_{y}(h)=0$ for $h>1$.

## Chapter 2

## Time Series Regression and EDA

### 2.1 Classical Regression for Time Series

We begin our discussion of linear regression in the time series context by assuming some output or dependent time series, say, $x_{t}$, for $t=1, \ldots, n$, is being influenced by a collection of possible inputs or independent series, say, $z_{t 1}, z_{t 2}, \ldots, z_{t q}$, where we first regard the inputs as fixed and known. This assumption, necessary for applying conventional linear regression, will be relaxed later on. We express this relation through the linear regression model

$$
\begin{equation*}
x_{t}=\beta_{0}+\beta_{1} z_{t 1}+\beta_{2} z_{t 2}+\cdots+\beta_{q} z_{t q}+w_{t} \tag{2.1}
\end{equation*}
$$

where $\beta_{0}, \beta_{1}, \ldots, \beta_{q}$ are unknown fixed regression coefficients, and $\left\{w_{t}\right\}$ is a random error or noise process consisting of independent and identically distributed (iid) normal variables with mean zero and variance $\sigma_{w}^{2}$; we will relax the iid assumption later.

## Example 2.1 Estimating a Linear Trend

Consider the global temperature data, say $x_{t}$, shown in Figure 1.3 and Figure 2.1. As discussed in Example 1.2, there is an apparent upward trend in the series that has been used to argue the global warming hypothesis. We might use simple linear regression to estimate that trend by fitting the model

$$
x_{t}=\beta_{0}+\beta_{1} z_{t}+w_{t}, \quad z_{t}=1880,1857, \ldots, 2009
$$

This is in the form of the regression model (2.1) with $q=1$. Note that we are making the assumption that the errors, $w_{t}$, are an iid normal sequence, which may not be true; the problem of autocorrelated errors is discussed in detail in Chapter 3. Also note that we could have used, for example, $z_{t}=1, \ldots, 130$, without affecting the interpretation of the slope coefficient, $\beta_{1}$; only the intercept, $\beta_{0}$, would change.

In ordinary least squares (OLS), we minimize the error sum of squares

$$
Q=\sum_{t=1}^{n} w_{t}^{2}=\sum_{t=1}^{n}\left(x_{t}-\left[\beta_{0}+\beta_{1} z_{t}\right]\right)^{2}
$$



Fig. 2.1. Global temperature deviations shown in Figure 1.3 with fitted linear trend line, which SUCKS as an estimate of trend.
with respect to $\beta_{i}$ for $i=0,1$. In this case we can use simple calculus to evaluate $\partial Q / \partial \beta_{i}=0$ for $i=0,1$, to obtain two equations to solve for the $\beta \mathrm{s}$. The OLS estimates of the coefficients are explicit and given by

$$
\widehat{\beta}_{1}=\frac{\sum_{t=1}^{n}\left(x_{t}-\bar{x}\right)\left(z_{t}-\bar{z}_{t}\right)}{\sum_{t=1}^{n}\left(z_{t}-\bar{z}\right)^{2}} \quad \text { and } \quad \widehat{\beta}_{0}=\bar{x}-\widehat{\beta}_{1} \bar{z}
$$

where $\bar{x}=\sum_{t} x_{t} / n$ and $\bar{z}_{t}=\sum_{t} z_{t} / n$ are the respective sample means. Note that we could have used, for example, $z_{t}=1, \ldots, 130$, without affecting the interpretation of the slope coefficient, $\beta_{1}$; only the intercept, $\beta_{0}$, would change.

Using R, we obtained the estimated coefficients $\widehat{\beta}_{0}=-11.2$, and $\widehat{\beta}_{1}=.006$ (with a standard error of .0003) yielding a highly significant estimated increase of .6 degrees centigrade per 100 years. We discuss the precise way in which the solution was accomplished after the example. Finally, Figure 2.1 shows the global temperature data, say $x_{t}$, with the estimated trend, say $\widehat{x}_{t}=-11.2+.006 z_{t}$, superimposed. It is apparent that the estimated trend line obtained via simple linear regression does not quite capture the trend of the data and better models will be needed. To perform this analysis in R, use the following commands:

```
summary(fit <- lm(gtemp~time(gtemp))) # regress gtemp on time - view results
plot(gtemp, type="o", ylab="Global Temperature Deviation")
abline(fit) # add the fitted regression line to the plot
```

The multiple linear regression model described by (2.1) can be conveniently written in a more general notation by defining the column vectors $z_{t}=\left(1, z_{t 1}, z_{t 2}, \ldots, z_{t q}\right)^{\prime}$ and $\boldsymbol{\beta}=\left(\beta_{0}, \beta_{1}, \ldots, \beta_{q}\right)^{\prime}$, where ' denotes transpose, so (2.1) can be written in the alternate form

$$
\begin{equation*}
x_{t}=\beta_{0}+\beta_{1} z_{t 1}+\cdots+\beta_{q} z_{t q}+w_{t}=\boldsymbol{\beta}^{\prime} z_{t}+w_{t} \tag{2.2}
\end{equation*}
$$

where $w_{t} \sim \operatorname{iid} \mathrm{~N}\left(0, \sigma_{w}^{2}\right)$. As in the previous example, OLS estimation finds the coefficient vector $\boldsymbol{\beta}$ that minimizes the error sum of squares

$$
\begin{equation*}
Q=\sum_{t=1}^{n} w_{t}^{2}=\sum_{t=1}^{n}\left(x_{t}-\boldsymbol{\beta}^{\prime} z_{t}\right)^{2} \tag{2.3}
\end{equation*}
$$

with respect to $\beta_{0}, \beta_{1}, \ldots, \beta_{q}$. This minimization can be accomplished by solving $\partial Q / \partial \beta_{i}=0$ for $i=0,1, \ldots, q$, which yields $q+1$ equations with $q+1$ unknowns. In vector notation, this procedure gives the normal equations

$$
\begin{equation*}
\left(\sum_{t=1}^{n} z_{t} z_{t}^{\prime}\right) \widehat{\boldsymbol{\beta}}=\sum_{t=1}^{n} z_{t} x_{t} \tag{2.4}
\end{equation*}
$$

If $\sum_{t=1}^{n} z_{t} z_{t}^{\prime}$ is non-singular, the least squares estimate of $\boldsymbol{\beta}$ is

$$
\widehat{\boldsymbol{\beta}}=\left(\sum_{t=1}^{n} z_{t} z_{t}^{\prime}\right)^{-1} \sum_{t=1}^{n} z_{t} x_{t}
$$

The minimized error sum of squares (2.3), denoted $S S E$, can be written as

$$
\begin{equation*}
S S E=\sum_{t=1}^{n}\left(x_{t}-\widehat{\boldsymbol{\beta}}^{\prime} z_{t}\right)^{2} \tag{2.5}
\end{equation*}
$$

The ordinary least squares estimators are unbiased, i.e., $E(\widehat{\boldsymbol{\beta}})=\boldsymbol{\beta}$, and have the smallest variance within the class of linear unbiased estimators.

If the errors $w_{t}$ are normally distributed, $\widehat{\boldsymbol{\beta}}$ is normally distributed with

$$
\begin{equation*}
\operatorname{cov}(\widehat{\boldsymbol{\beta}})=\sigma_{w}^{2} C \tag{2.6}
\end{equation*}
$$

where

$$
\begin{equation*}
C=\left(\sum_{t=1}^{n} z_{t} z_{t}^{\prime}\right)^{-1} \tag{2.7}
\end{equation*}
$$

is a convenient notation. An unbiased estimator for the variance $\sigma_{w}^{2}$ is

$$
\begin{equation*}
s_{w}^{2}=M S E=\frac{S S E}{n-(q+1)}, \tag{2.8}
\end{equation*}
$$

where MSE denotes the mean squared error. Under the normal assumption,

$$
\begin{equation*}
t=\frac{\left(\widehat{\beta}_{i}-\beta_{i}\right)}{s_{w} \sqrt{c_{i i}}} \tag{2.9}
\end{equation*}
$$

has the t -distribution with $n-(q+1)$ degrees of freedom; $c_{i i}$ denotes the $i$-th diagonal element of $C$, as defined in (2.7). This result is often used for individual tests of the null hypothesis $\mathrm{H}_{0}: \beta_{i}=0$ for $i=1, \ldots, q$.

Various competing models are often of interest to isolate or select the best subset of independent variables. Suppose a proposed model specifies that only a subset $r<q$ independent variables, say, $z_{t, 1: r}=\left\{z_{t 1}, z_{t 2}, \ldots, z_{t r}\right\}$ is influencing the dependent variable $x_{t}$. The reduced model is

Table 2.1. Analysis of Variance for Regression

| Source | df | Sum of Squares | Mean Square | $F$ |
| :--- | :---: | :---: | :---: | :---: |
| $z_{t, r+1: q}$ | $q-r$ | $S S R=S S E_{r}-S S E$ | $M S R=\operatorname{SSR} /(q-r)$ | $F=\frac{M S R}{M S E}$ |
| Error | $n-(q+1)$ | $S S E$ | $M S E=\operatorname{SSE} /(n-q-1)$ |  |

$$
\begin{equation*}
x_{t}=\beta_{0}+\beta_{1} z_{t 1}+\cdots+\beta_{r} z_{t r}+w_{t} \tag{2.10}
\end{equation*}
$$

where $\beta_{1}, \beta_{2}, \ldots, \beta_{r}$ are a subset of coefficients of the original $q$ variables.
The null hypothesis in this case is $\mathrm{H}_{0}: \beta_{r+1}=\cdots=\beta_{q}=0$. We can test the reduced model (2.10) against the full model (2.2) by comparing the error sums of squares under the two models using the $F$-statistic

$$
\begin{equation*}
F=\frac{\left(S S E_{r}-S S E\right) /(q-r)}{S S E /(n-q-1)}=\frac{M S R}{M S E} \tag{2.11}
\end{equation*}
$$

where $\operatorname{SSE} E_{r}$ is the error sum of squares under the reduced model (2.10). Note that $S S E_{r} \geq S S E$ because the full model has more parameters. If $\mathrm{H}_{0}: \beta_{r+1}=\cdots=\beta_{q}=$ 0 is true, then $S S E_{r} \approx S S E$ because the estimates of those $\beta$ s will be close to 0 . Hence, we do not believe $\mathrm{H}_{0}$ if $S S R=S S E_{r}-S S E$ is big. Under the null hypothesis, (2.11) has a central $F$-distribution with $q-r$ and $n-q-1$ degrees of freedom when (2.10) is the correct model.

These results are often summarized in an Analysis of Variance (ANOVA) table as given in Table 2.1 for this particular case. The difference in the numerator is often called the regression sum of squares ( $S S R$ ). The null hypothesis is rejected at level $\alpha$ if $F>F_{n-q-1}^{q-r}(\alpha)$, the $1-\alpha$ percentile of the $F$ distribution with $q-r$ numerator and $n-q-1$ denominator degrees of freedom.

A special case of interest is the null hypothesis $\mathrm{H}_{0}: \beta_{1}=\cdots=\beta_{q}=0$. In this case $r=0$, and the model in (2.10) becomes

$$
x_{t}=\beta_{0}+w_{t} .
$$

We may measure the proportion of variation accounted for by all the variables using

$$
\begin{equation*}
R^{2}=\frac{S S E_{0}-S S E}{S S E_{0}}, \tag{2.12}
\end{equation*}
$$

where the residual sum of squares under the reduced model is

$$
\begin{equation*}
S S E_{0}=\sum_{t=1}^{n}\left(x_{t}-\bar{x}\right)^{2} . \tag{2.13}
\end{equation*}
$$

In this case $S S E_{0}$ is the sum of squared deviations from the mean $\bar{x}$ and is otherwise known as the adjusted total sum of squares. The measure $R^{2}$ is called the squared multiple correlation between $x_{t}$ and the variables $z_{t 1}, z_{t 2}, \ldots, z_{t q}$.

The techniques discussed in the previous paragraph can be used to test various models against one another using the $F$ test given in (2.11). These tests have been used in the past in a stepwise manner, where variables are added or deleted when the values from the $F$-test either exceed or fail to exceed some predetermined levels. The procedure, called stepwise multiple regression, is useful in arriving at a set of useful variables. An alternative is to focus on a procedure for model selection that does not proceed sequentially, but simply evaluates each model on its own merits. Suppose we consider a normal regression model with $k$ coefficients and denote the maximum likelihood estimator for the variance as

$$
\begin{equation*}
\widehat{\sigma}_{k}^{2}=\frac{\operatorname{SSE}(k)}{n}, \tag{2.14}
\end{equation*}
$$

where $\operatorname{SSE}(k)$ denotes the residual sum of squares under the model with $k$ regression coefficients. Then, Akaike $(1969,1973,1974)$ suggested measuring the goodness of fit for this particular model by balancing the error of the fit against the number of parameters in the model; we define the following.

## Definition 2.1 Akaike's Information Criterion (AIC)

$$
\begin{equation*}
\mathrm{AIC}=\log \widehat{\sigma}_{k}^{2}+\frac{n+2 k}{n} \tag{2.15}
\end{equation*}
$$

where $\widehat{\sigma}_{k}^{2}$ is given by (2.14) and $k$ is the number of parameters in the model.
The value of $k$ yielding the minimum AIC specifies the best model. ${ }^{1}$ The idea is roughly that minimizing $\widehat{\sigma}_{k}^{2}$ would be a reasonable objective, except that it decreases monotonically as $k$ increases. Therefore, we ought to penalize the error variance by a term proportional to the number of parameters. The choice for the penalty term given by (2.15) is not the only one, and a considerable literature is available advocating different penalty terms. A corrected form, suggested by Sugiura (1978), and expanded by Hurvich and Tsai (1989), can be based on small-sample distributional results for the linear regression model. The corrected form is defined as follows.

Definition 2.2 AIC, Bias Corrected (AICc)

$$
\begin{equation*}
\mathrm{AICc}=\log \widehat{\sigma}_{k}^{2}+\frac{n+k}{n-k-2} \tag{2.16}
\end{equation*}
$$

where $\widehat{\sigma}_{k}^{2}$ is given by (2.14), $k$ is the number of parameters in the model, and $n$ is the sample size.

We may also derive a correction term based on Bayesian arguments, as in Schwarz (1978), which leads to the following.

[^4]
## Definition 2.3 Bayesian Information Criterion (BIC)

$$
\begin{equation*}
\mathrm{BIC}=\log \widehat{\sigma}_{k}^{2}+\frac{k \log n}{n}, \tag{2.17}
\end{equation*}
$$

using the same notation as in Definition 2.2.
BIC is also called the Schwarz Information Criterion (SIC); see also Rissanen (1978) for an approach yielding the same statistic based on a minimum description length argument. Various simulation studies have tended to verify that BIC does well at getting the correct order in large samples, whereas AICc tends to be superior in smaller samples where the relative number of parameters is large; see McQuarrie and Tsai (1998) for detailed comparisons. In fitting regression models, two measures that have been used in the past are adjusted R-squared, which is essentially $s_{w}^{2}$, and Mallows $C_{p}$, Mallows (1973), which we do not consider in this context.

## Example 2.2 Pollution, Temperature and Mortality

The data shown in Figure 2.2 are extracted series from a study by Shumway et al. (1988) of the possible effects of temperature and pollution on weekly mortality in Los Angeles County. Note the strong seasonal components in all of the series, corresponding to winter-summer variations and the downward trend in the cardiovascular mortality over the 10-year period.

A scatterplot matrix, shown in Figure 2.3, indicates a possible linear relation between mortality and the pollutant particulates and a possible relation to temperature. Note the curvilinear shape of the temperature mortality curve, indicating that higher temperatures as well as lower temperatures are associated with increases in cardiovascular mortality.

Based on the scatterplot matrix, we entertain, tentatively, four models where $M_{t}$ denotes cardiovascular mortality, $T_{t}$ denotes temperature and $P_{t}$ denotes the particulate levels. They are

$$
\begin{align*}
& M_{t}=\beta_{1}+\beta_{2} t+w_{t}  \tag{2.18}\\
& M_{t}=\beta_{1}+\beta_{2} t+\beta_{3}\left(T_{t}-T .\right)+w_{t}  \tag{2.19}\\
& M_{t}=\beta_{1}+\beta_{2} t+\beta_{3}\left(T_{t}-T .\right)+\beta_{4}\left(T_{t}-T .\right)^{2}+w_{t}  \tag{2.20}\\
& M_{t}=\beta_{1}+\beta_{2} t+\beta_{3}\left(T_{t}-T .\right)+\beta_{4}\left(T_{t}-T .\right)^{2}+\beta_{5} P_{t}+w_{t} \tag{2.21}
\end{align*}
$$

where we adjust temperature for its mean, $T .=74.6$, to avoid scaling problems. It is clear that (2.18) is a trend only model, (2.19) is linear temperature, (2.20) is curvilinear temperature and (2.21) is curvilinear temperature and pollution. We summarize some of the statistics given for this particular case in Table 2.2. The values of $R^{2}$ were computed by noting that $S S E_{1}=50,687$ using (2.13).

We note that each model does substantially better than the one before it and that the model including temperature, temperature squared, and particulates does the best, accounting for some $60 \%$ of the variability and with the best value for AIC and BIC (because of the large sample size, AIC and AICc are nearly the same). Note that one can compare any two models using the residual sums of squares and


Fig. 2.2. Average weekly cardiovascular mortality (top), temperature (middle) and particulate pollution (bottom) in Los Angeles County. There are 508 six-day smoothed averages obtained by filtering daily values over the 10 year period 1970-1979.
(2.11). Hence, a model with only trend could be compared to the full model using $q=5, r=2, n=508$, so

$$
F_{3,503}=\frac{(40,020-20,508) / 3}{20,508 / 503}=160,
$$

which exceeds $F_{3,503}(.001)=5.51$. We obtain the best prediction model,

$$
\begin{aligned}
\widehat{M}_{t}= & 81.59-.027_{(.002)} t-.473_{(.032)}\left(T_{t}-74.6\right) \\
& +.023_{(.003)}\left(T_{t}-74.6\right)^{2}+.255_{(.019)} P_{t}
\end{aligned}
$$

for mortality, where the standard errors, computed from (2.6)-(2.8), are given in parentheses. As expected, a negative trend is present in time as well as a negative coefficient for adjusted temperature. The quadratic effect of temperature can clearly be seen in the scatterplots of Figure 2.3. Pollution weights positively and can be interpreted as the incremental contribution to daily deaths per unit of particulate pollution. It would still be essential to check the residuals $\widehat{w}_{t}=M_{t}-\widehat{M}_{t}$ for autocorrelation (of which there is a substantial amount), but we defer this question to to $\S 3.9$ when we discuss regression with correlated errors.

Below is the R code to plot the series, display the scatterplot matrix, fit the final regression model (2.21), and compute the corresponding values of AIC and BIC.


Fig. 2.3. Scatterplot matrix showing plausible relations between mortality, temperature, and pollution.

Table 2.2. Summary Statistics for Mortality Models

| Model | $k$ | SSE | df | MSE | $R^{2}$ | AIC | BIC |
| :---: | :---: | ---: | :---: | :---: | :---: | :---: | :---: |
| $(2.18)$ | 2 | 40,020 | 506 | 79.0 | .21 | 5.38 | 5.40 |
| $(2.19)$ | 3 | 31,413 | 505 | 62.2 | .38 | 5.14 | 5.17 |
| $(2.20)$ | 4 | 27,985 | 504 | 55.5 | .45 | 5.03 | 5.07 |
| $(2.21)$ | 5 | 20,508 | 503 | 40.8 | .60 | 4.72 | 4.77 |

Finally, the use of na.action in $\operatorname{lm}()$ is to retain the time series attributes for the residuals and fitted values.

```
par(mfrow=c (3,1))
plot(cmort, main="Cardiovascular Mortality", xlab="", ylab="")
plot(tempr, main="Temperature", xlab="", ylab="")
plot(part, main="Particulates", xlab="", ylab="")
dev.new() # open a new graphic device for the scatterplot matrix
pairs(cbind(Mortality=cmort, Temperature=tempr, Particulates=part))
temp = tempr-mean(tempr) # center temperature
temp2 = temp^2
trend = time(cmort) # time
fit = lm(cmort~ trend + temp + temp2 + part, na.action=NULL)
summary(fit) # regression results
summary(aov(fit)) # ANOVA table (compare to next line)
summary(aov(lm(cmort~cbind(trend, temp, temp2, part)))) # Table 2.1
num = length(cmort) # sample size
```

```
AIC(fit)/num - log(2*pi) # AIC as in Def 2.1
BIC(fit)/num - log(2*pi) # BIC as in Def 2.3
```

As previously mentioned, it is possible to include lagged variables in time series regression models and we will continue to discuss this type of problem throughout the text. This concept is explored further in Problem 2.2. The following is a simple example of lagged regression.

## Example 2.3 Regression With Lagged Variables

In Example 1.25, we discovered that the Southern Oscillation Index (SOI) measured at time $t-6$ months is associated with the Recruitment series at time $t$, indicating that the SOI leads the Recruitment series by six months. Although there is strong evidence that the relationship is NOT linear (this is discussed further in Example 2.7), for demonstration purposes only, we consider the following regression,

$$
\begin{equation*}
R_{t}=\beta_{1}+\beta_{2} S_{t-6}+w_{t}, \tag{2.22}
\end{equation*}
$$

where $R_{t}$ denotes Recruitment for month $t$ and $S_{t-6}$ denotes SOI six months prior. Assuming the $w_{t}$ sequence is white, the fitted model is

$$
\begin{equation*}
\widehat{R}_{t}=65.79-44.28_{(2.78)} S_{t-6} \tag{2.23}
\end{equation*}
$$

with $\widehat{\sigma}_{w}=22.5$ on 445 degrees of freedom. This result indicates the strong predictive ability of SOI for Recruitment six months in advance. Of course, it is still essential to check the the model assumptions, but we defer this discussion until later.

Performing lagged regression in R is a little difficult because the series must be aligned prior to running the regression. The easiest way to do this is to create an object that we call fish using ts.intersect, which aligns the lagged series.
fish = ts.intersect ( rec, soiL6=lag(soi,-6) )
summary (fit1 <- lm(rec~ soiL6, data=fish, na.action=NULL))
The headache of aligning the lagged series can be avoided by using the R package dynlm, which must be downloaded and installed.
require (dynlm)
summary (fit2 <- dynlm(rec~ L(soi,6)))
In the dynlm example, fit2 is similar to a 1 m object, but the time series attributes are retained without any additional commands.

### 2.2 Exploratory Data Analysis

In general, it is necessary for time series data to be stationary so averaging lagged products over time, as in the previous section, will be a sensible thing to do. With time series data, it is the dependence between the values of the series that is important to measure; we must, at least, be able to estimate autocorrelations with precision. It would be difficult to measure that dependence if the dependence structure is not
regular or is changing at every time point. Hence, to achieve any meaningful statistical analysis of time series data, it will be crucial that, if nothing else, the mean and the autocovariance functions satisfy the conditions of stationarity (for at least some reasonable stretch of time) stated in Definition 1.7. Often, this is not the case, and we will mention some methods in this section for playing down the effects of nonstationarity so the stationary properties of the series may be studied.

A number of our examples came from clearly nonstationary series. The Johnson \& Johnson series in Figure 1.1 has a mean that increases exponentially over time, and the increase in the magnitude of the fluctuations around this trend causes changes in the covariance function; the variance of the process, for example, clearly increases as one progresses over the length of the series. Also, the global temperature series shown in Figure 1.3 contains some evidence of a trend over time; human-induced global warming advocates seize on this as empirical evidence to advance their hypothesis that temperatures are increasing.

Perhaps the easiest form of nonstationarity to work with is the trend stationary model wherein the process has stationary behavior around a trend. We may write this type of model as

$$
\begin{equation*}
x_{t}=\mu_{t}+y_{t} \tag{2.24}
\end{equation*}
$$

where $x_{t}$ are the observations, $\mu_{t}$ denotes the trend, and $y_{t}$ is a stationary process. Quite often, strong trend, $\mu_{t}$, will obscure the behavior of the stationary process, $y_{t}$, as we shall see in numerous examples. Hence, there is some advantage to removing the trend as a first step in an exploratory analysis of such time series. The steps involved are to obtain a reasonable estimate of the trend component, say $\widehat{\mu}_{t}$, and then work with the residuals

$$
\begin{equation*}
\widehat{y}_{t}=x_{t}-\widehat{\mu}_{t} . \tag{2.25}
\end{equation*}
$$

Consider the following example.

## Example 2.4 Detrending Global Temperature

Here we suppose the model is of the form of (2.24),

$$
x_{t}=\mu_{t}+y_{t},
$$

where, as we suggested in the analysis of the global temperature data presented in Example 2.1, a straight line might be a reasonable model for the trend, i.e.,

$$
\mu_{t}=\beta_{1}+\beta_{2} t
$$

In that example, we estimated the trend using ordinary least squares ${ }^{2}$ and found

$$
\widehat{\mu}_{t}=-11.2+.006 t
$$

[^5]

Fig. 2.4. Detrended (top) and differenced (bottom) global temperature series. The original data are shown in Figure 1.3 and Figure 2.1.

Figure 2.1 shows the data with the estimated trend line superimposed. To obtain the detrended series we simply subtract $\widehat{\mu}_{t}$ from the observations, $x_{t}$, to obtain the detrended series

$$
\widehat{y}_{t}=x_{t}+11.2-.006 t .
$$

The top graph of Figure 2.4 shows the detrended series. Figure 2.5 shows the ACF of the original data (top panel) as well as the ACF of the detrended data (middle panel).

To detrend in the series in R, use the following commands. We also show how to difference and plot the differenced data; we discuss differencing after this example. In addition, we show how to generate the sample ACFs displayed in Figure 2.5.

```
fit = lm(gtemp~time(gtemp), na.action=NULL) # regress gtemp on time
par(mfrow=c(2,1))
plot(resid(fit), type="o", main="detrended")
plot(diff(gtemp), type="o", main="first difference")
par(mfrow=c(3,1)) # plot ACFs
acf(gtemp, 48, main="gtemp")
acf(resid(fit), 48, main="detrended")
acf(diff(gtemp), 48, main="first difference")
```

In Example 1.9 and the corresponding Figure 1.9 we saw that a random walk might also be a good model for trend. That is, rather than modeling trend as fixed (as in Example 2.4), we might model trend as a stochastic component using the random walk with drift model,

$$
\begin{equation*}
\mu_{t}=\delta+\mu_{t-1}+w_{t}, \tag{2.26}
\end{equation*}
$$

where $w_{t}$ is white noise and is independent of $y_{t}$. If the appropriate model is (2.24), then differencing the data, $x_{t}$, yields a stationary process; that is,

$$
\begin{align*}
x_{t}-x_{t-1} & =\left(\mu_{t}+y_{t}\right)-\left(\mu_{t-1}+y_{t-1}\right)  \tag{2.27}\\
& =\delta+w_{t}+y_{t}-y_{t-1} .
\end{align*}
$$

It is easy to show $z_{t}=y_{t}-y_{t-1}$ is stationary using Property 1.1 of Chapter 1. That is, because $y_{t}$ is stationary,

$$
\begin{align*}
\gamma_{z}(h) & =\operatorname{cov}\left(z_{t+h}, z_{t}\right)=\operatorname{cov}\left(y_{t+h}-y_{t+h-1}, y_{t}-y_{t-1}\right) \\
& =2 \gamma_{y}(h)-\gamma_{y}(h+1)-\gamma_{y}(h-1) \tag{2.28}
\end{align*}
$$

is independent of time; we leave it as an exercise (Problem 2.5) to show that $x_{t}-x_{t-1}$ in (2.27) is stationary.

One advantage of differencing over detrending to remove trend is that no parameters are estimated in the differencing operation. One disadvantage, however, is that differencing does not yield an estimate of the stationary process $y_{t}$ as can be seen in (2.27). If an estimate of $y_{t}$ is essential, then detrending may be more appropriate. If the goal is to coerce the data to stationarity, then differencing may be more appropriate. Differencing is also a viable tool if the trend is fixed, as in Example 2.4. That is, e.g., if $\mu_{t}=\beta_{1}+\beta_{2} t$ in the model (2.24), differencing the data produces stationarity (see Problem 2.4):

$$
x_{t}-x_{t-1}=\left(\mu_{t}+y_{t}\right)-\left(\mu_{t-1}+y_{t-1}\right)=\beta_{2}+y_{t}-y_{t-1} .
$$

Because differencing plays a central role in time series analysis, it receives its own notation. The first difference is denoted as

$$
\begin{equation*}
\nabla x_{t}=x_{t}-x_{t-1} . \tag{2.29}
\end{equation*}
$$

As we have seen, the first difference eliminates a linear trend. A second difference, that is, the difference of (2.29), can eliminate a quadratic trend, and so on. In order to define higher differences, we need a variation in notation that we will use often in our discussion of ARIMA models in Chapter 3.

Definition 2.4 We define the backshift operator by

$$
B x_{t}=x_{t-1}
$$

and extend it to powers $B^{2} x_{t}=B\left(B x_{t}\right)=B x_{t-1}=x_{t-2}$, and so on. Thus,

$$
\begin{equation*}
B^{k} x_{t}=x_{t-k} \tag{2.30}
\end{equation*}
$$

The idea of an inverse operator can also be given if we require $B^{-1} B=1$, so that

$$
x_{t}=B^{-1} B x_{t}=B^{-1} x_{t-1}
$$

That is, $B^{-1}$ is the forward-shift operator. In addition, it is clear that we may rewrite (2.29) as


Fig. 2.5. Sample ACFs of the global temperature (top), and of the detrended (middle) and the differenced (bottom) series.

$$
\begin{equation*}
\nabla x_{t}=(1-B) x_{t}, \tag{2.31}
\end{equation*}
$$

and we may extend the notion further. For example, the second difference becomes

$$
\begin{equation*}
\nabla^{2} x_{t}=(1-B)^{2} x_{t}=\left(1-2 B+B^{2}\right) x_{t}=x_{t}-2 x_{t-1}+x_{t-2} \tag{2.32}
\end{equation*}
$$

by the linearity of the operator. To check, just take the difference of the first difference $\nabla\left(\nabla x_{t}\right)=\nabla\left(x_{t}-x_{t-1}\right)=\left(x_{t}-x_{t-1}\right)-\left(x_{t-1}-x_{t-2}\right)$.

Definition 2.5 Differences of order $\mathbf{d}$ are defined as

$$
\begin{equation*}
\nabla^{d}=(1-B)^{d} \tag{2.33}
\end{equation*}
$$

where we may expand the operator $(1-B)^{d}$ algebraically to evaluate for higher integer values of $d$. When $d=1$, we drop it from the notation.

The first difference (2.29) is an example of a linear filter applied to eliminate a trend. Other filters, formed by averaging values near $x_{t}$, can produce adjusted series that eliminate other kinds of unwanted fluctuations, as in Chapter 3. The differencing technique is an important component of the ARIMA model discussed in Chapter 3.

## Example 2.5 Differencing Global Temperature

The first difference of the global temperature series, also shown in Figure 2.4, produces different results than removing trend by detrending via regression. For example, the differenced series does not contain the long middle cycle we observe in the detrended series. The ACF of this series is also shown in Figure 2.5. In this case it appears that the differenced process shows minimal autocorrelation, which may imply the global temperature series is nearly a random walk with drift. It is interesting to note that if the series is a random walk with drift, the mean of the differenced series, which is an estimate of the drift, is about .0066 (but with a large standard error):
mean(diff(gtemp)) $\quad=0.00659$ (drift)
sd(diff(gtemp))/sqrt(length(diff(gtemp))) \# = 0.00966 (SE)
Often, obvious aberrations are present that can contribute nonstationary as well as nonlinear behavior in observed time series. In such cases, transformations may be useful to equalize the variability over the length of a single series. A particularly useful transformation is

$$
\begin{equation*}
y_{t}=\log x_{t}, \tag{2.34}
\end{equation*}
$$

which tends to suppress larger fluctuations that occur over portions of the series where the underlying values are larger. Other possibilities are power transformations in the Box-Cox family of the form

$$
y_{t}= \begin{cases}\left(x_{t}^{\lambda}-1\right) / \lambda & \lambda \neq 0  \tag{2.35}\\ \log x_{t} & \lambda=0\end{cases}
$$

Methods for choosing the power $\lambda$ are available (see Johnson and Wichern, 1992, §4.7) but we do not pursue them here. Often, transformations are also used to improve the approximation to normality or to improve linearity in predicting the value of one series from another.

## Example 2.6 Paleoclimatic Glacial Varves

Melting glaciers deposit yearly layers of sand and silt during the spring melting seasons, which can be reconstructed yearly over a period ranging from the time deglaciation began in New England (about 12,600 years ago) to the time it ended (about 6,000 years ago). Such sedimentary deposits, called varves, can be used as proxies for paleoclimatic parameters, such as temperature, because, in a warm year, more sand and silt are deposited from the receding glacier. Figure 2.6 shows the thicknesses of the yearly varves collected from one location in Massachusetts for 634 years, beginning 11,834 years ago. For further information, see Shumway and Verosub (1992). Because the variation in thicknesses increases in proportion to the amount deposited, a logarithmic transformation could remove the nonstationarity observable in the variance as a function of time. Figure 2.6 shows the original and transformed varves, and it is clear that this improvement has occurred. We may also plot the histogram of the original and transformed data, as in Problem 2.6, to argue that the approximation to normality is improved. The ordinary first differences (2.31) are also computed in Problem 2.6, and we note that the first differences have


Fig. 2.6. Glacial varve thicknesses (top) from Massachusetts for $n=634$ years compared with log transformed thicknesses (bottom).
a significant negative correlation at lag $h=1$. Later, in Chapter 5, we will show that perhaps the varve series has long memory and will propose using fractional differencing.

Figure 2.6 was generated in R as follows:

```
par(mfrow=c (2,1))
plot(varve, main="varve", ylab="")
plot(log(varve), main="log(varve)", ylab="" )
```

Next, we consider another preliminary data processing technique that is used for the purpose of visualizing the relations between series at different lags, namely, scatterplot matrices. In the definition of the ACF, we are essentially interested in relations between $x_{t}$ and $x_{t-h}$; the autocorrelation function tells us whether a substantial linear relation exists between the series and its own lagged values. The ACF gives a profile of the linear correlation at all possible lags and shows which values of $h$ lead to the best predictability. The restriction of this idea to linear predictability, however, may mask a possible nonlinear relation between current values, $x_{t}$, and past values, $x_{t-h}$. This idea extends to two series where one may be interested in examining scatterplots of $y_{t}$ versus $x_{t-h}$.

## Example 2.7 Scatterplot Matrices, SOI and Recruitment

To check for nonlinear relations of this form, it is convenient to display a lagged scatterplot matrix, as in Figure 2.7, that displays values of the SOI, $S_{t}$, on the vertical axis plotted against $S_{t-h}$ on the horizontal axis. The sample autocorrelations are displayed in the upper right-hand corner and superimposed on the scatterplots are locally weighted scatterplot smoothing (lowess) lines that can be used to help


Fig. 2.7. Scatterplot matrix relating current SOI values, $S_{t}$, to past SOI values, $S_{t-h}$, at lags $h=1,2, \ldots, 12$. The values in the upper right corner are the sample autocorrelations and the lines are a lowess fit.
discover any nonlinearities. We discuss smoothing in the next section, but for now, think of lowess as a robust method for fitting nonlinear regression.

In Figure 2.7, we notice that the lowess fits are approximately linear, so that the sample autocorrelations are meaningful. Also, we see strong positive linear relations at lags $h=1,2,11,12$, that is, between $S_{t}$ and $S_{t-1}, S_{t-2}, S_{t-11}, S_{t-12}$, and a negative linear relation at lags $h=6,7$. These results match up well with peaks noticed in the ACF in Figure 1.16.

Similarly, we might want to look at values of one series, say Recruitment, denoted $R_{t}$ plotted against another series at various lags, say the SOI, $S_{t-h}$, to look for possible nonlinear relations between the two series. Because, for example, we might wish to predict the Recruitment series, $R_{t}$, from current or past values of the SOI series, $S_{t-h}$, for $h=0,1,2, \ldots$ it would be worthwhile to examine the scatterplot matrix. Figure 2.8 shows the lagged scatterplot of the Recruitment series $R_{t}$ on the


Fig. 2.8. Scatterplot matrix of the Recruitment series, $R_{t}$, on the vertical axis plotted against the SOI series, $S_{t-h}$, on the horizontal axis at lags $h=0,1, \ldots, 8$. The values in the upper right corner are the sample cross-correlations and the lines are a lowess fit.
vertical axis plotted against the SOI index $S_{t-h}$ on the horizontal axis. In addition, the figure exhibits the sample cross-correlations as well as lowess fits.

Figure 2.8 shows a fairly strong nonlinear relationship between Recruitment, $R_{t}$, and the SOI series at $S_{t-5}, S_{t-6}, S_{t-7}, S_{t-8}$, indicating the SOI series tends to lead the Recruitment series and the coefficients are negative, implying that increases in the SOI lead to decreases in the Recruitment. The nonlinearity observed in the scatterplots (with the help of the superimposed lowess fits) indicate that the behavior between Recruitment and the SOI is different for positive values of SOI than for negative values of SOI.

Simple scatterplot matrices for one series can be obtained in R using the lag.plot command. Figure 2.7 and Figure 2.8 may be reproduced using the following scripts provided with the text (see Appendix R for detials):

```
lag1.plot(soi, 12) # Figure 2.7
lag2.plot(soi, rec, 8) # Figure 2.8
```



Fig. 2.9. Display for Example 2.8: Plot of Recruitment $\left(R_{t}\right)$ vs SOI lagged 6 months ( $S_{t-6}$ ) with the fitted values of the regression as points ( + ) and a lowess fit (一).

## Example 2.8 Regression with Lagged Variables (cont)

In Example 2.3 we regressed Recruitment on lagged SOI,

$$
R_{t}=\beta_{1}+\beta_{2} S_{t-6}+w_{t} .
$$

However, in Example 2.7, we saw that the relationship is nonlinear and different when SOI is positive or negative. In this case, we may consider adding a dummy variable to account for this change. In particular, we fit the model

$$
R_{t}=\beta_{1}+\beta_{2} S_{t-6}+\beta_{3} D_{t-6}+\beta_{4} D_{t-6} S_{t-6}+w_{t},
$$

where $D_{t}$ is a dummy variable that is 0 if $S_{t}<0$ and 1 otherwise. The result of the fit is given in the R code below. Figure 2.9 shows $R_{t}$ vs $S_{t-6}$ with the fitted values of the regression and a lowess fit superimposed. The piecewise regression fit is similar to the lowess fit, but we note that the residuals are not white noise (see the code below). This is followed up in Example 3.40.
dummy = ifelse(soi<0, 0, 1)
fish = ts.intersect(rec, soiL6=lag(soi,-6), dL6=lag(dummy,-6), dframe=TRUE)
summary (fit <- lm(rec~ soiL6*dL6, data=fish, na.action=NULL))

Coefficients:

|  | Estimate | Std.Error | t.value | p.value |
| :--- | ---: | ---: | ---: | ---: |
| (Intercept) | 74.479 | 2.865 | 25.998 | $<2 \mathrm{e}-16$ |
| soiL6 | -15.358 | 7.401 | -2.075 | 0.0386 |
| dL6 | -1.139 | 3.711 | -0.307 | 0.7590 |
| soiL6:dL6 | -51.244 | 9.523 | -5.381 | $1.2 \mathrm{e}-07$ |

## --

Residual standard error: 21.84 on 443 degrees of freedom
Multiple R-squared: 0.4024
F-statistic: 99.43 on 3 and 443 DF, p-value: < 2.2e-16
attach (fish)
plot(soiL6, rec)
lines(lowess(soiL6, rec), col=4, lwd=2)
points(soiL6, fitted(fit), pch='+', col=2)
plot(resid(fit)) \# not shown ...
acf(resid(fit)) \# ... but obviously not noise


Fig. 2.10. Data generated by (2.36) [top] and the fitted line (2.38) superimposed on the data [bottom].

As a final exploratory tool, we discuss assessing periodic behavior in time series data using regression analysis and the periodogram; this material may be thought of as an introduction to spectral analysis, which we discuss in detail in Chapter 4. In Example 1.10, we briefly discussed the problem of identifying cyclic or periodic signals in time series. A number of the time series we have seen so far exhibit periodic behavior. For example, the data from the pollution study example shown in Figure 2.2 exhibit strong yearly cycles. Also, the Johnson \& Johnson data shown in Figure 1.1 make one cycle every year (four quarters) on top of an increasing trend and the speech data in Figure 1.3 is highly repetitive. The monthly SOI and Recruitment series in Figure 1.6 show strong yearly cycles, but hidden in the series are clues to the El Niño cycle.

## Example 2.9 Using Regression to Discover a Signal in Noise

In Example 1.10, we generated $n=500$ observations from the model

$$
\begin{equation*}
x_{t}=A \cos (2 \pi \omega t+\phi)+w_{t}, \tag{2.36}
\end{equation*}
$$

where $\omega=1 / 50, A=2, \phi=.6 \pi$, and $\sigma_{w}=5$; the data are shown on the bottom panel of Figure 1.10 on page 10. At this point we assume the frequency of oscillation $\omega=1 / 50$ is known, but $A$ and $\phi$ are unknown parameters. In this case the parameters appear in (2.36) in a nonlinear way, so we use a trigonometric identity ${ }^{3}$ and write

$$
A \cos (2 \pi \omega t+\phi)=\beta_{1} \cos (2 \pi \omega t)+\beta_{2} \sin (2 \pi \omega t)
$$

where $\beta_{1}=A \cos (\phi)$ and $\beta_{2}=-A \sin (\phi)$. Now the model (2.36) can be written in the usual linear regression form given by (no intercept term is needed here)

[^6]\[

$$
\begin{equation*}
x_{t}=\beta_{1} \cos (2 \pi t / 50)+\beta_{2} \sin (2 \pi t / 50)+w_{t} . \tag{2.37}
\end{equation*}
$$

\]

Using linear regression on the generated data, the fitted model is

$$
\begin{equation*}
\widehat{x}_{t}=-.71_{(.30)} \cos (2 \pi t / 50)-2.55_{(.30)} \sin (2 \pi t / 50) \tag{2.38}
\end{equation*}
$$

with $\widehat{\sigma}_{w}=4.68$, where the values in parentheses are the standard errors. We note the actual values of the coefficients for this example are $\beta_{1}=2 \cos (.6 \pi)=-.62$ and $\beta_{2}=-2 \sin (.6 \pi)=-1.90$. Because the parameter estimates are significant and close to the actual values, it is clear that we are able to detect the signal in the noise using regression, even though the signal appears to be obscured by the noise in the bottom panel of Figure 1.10. Figure 2.10 shows data generated by (2.36) with the fitted line, (2.38), superimposed.

To reproduce the analysis and Figure 2.10 in R, use the following:
set.seed(1000) \# so you can reproduce these results
$\mathrm{x}=2 * \cos (2 * \mathrm{pi} * 1: 500 / 50+.6 * \mathrm{pi})+\operatorname{rnorm}(500,0,5)$
$\mathrm{z} 1=\cos (2 * \mathrm{pi} * 1: 500 / 50)$
$z 2=\sin (2 * \mathrm{pi} * 1: 500 / 50)$
summary (fit <- $\operatorname{lm}(\mathrm{x} \sim 0+\mathrm{z} 1+\mathrm{z} 2))$ \# zero to exclude the intercept
$\operatorname{par}(\operatorname{mfrow}=c(2,1), \operatorname{mar}=c(3,3,1,1), \operatorname{mgp}=c(1.6, .6,0))$
plot.ts(x)
plot.ts(x, col=8, ylab=expression(hat(x)))
lines(fitted(fit), col=2)
We will discuss this and related approaches in more detail in Chapter 4.

### 2.3 Smoothing Time Series

In §1.4, we introduced the concept of smoothing a time series, and in Example 1.7, we discussed using a moving average to smooth white noise. This method is useful for discovering certain traits in a time series, such as long-term trend and seasonal components (see Section 4.7 for details). In particular, if $x_{t}$ represents the observations, then

$$
\begin{equation*}
m_{t}=\sum_{j=-k}^{k} a_{j} x_{t-j} \tag{2.39}
\end{equation*}
$$

where $a_{j}=a_{-j} \geq 0$ and $\sum_{j=-k}^{k} a_{j}=1$ is a symmetric moving average of the data.

## Example 2.10 Moving Average Smoother

For example, Figure 2.11 shows the monthly SOI series discussed in Example 1.4 smoothed using (2.39) with weights $a_{0}=a_{ \pm 1}=\cdots=a_{ \pm 5}=1 / 12$, and $a_{ \pm 6}=1 / 24$; $k=6$. This particular method removes (filters out) the obvious annual temperature cycle and helps emphasize the El Niño cycle.

To reproduce Figure 2.11 in R:

```
wgts = c(.5, rep (1,11), .5)/12
soif = filter(soi, sides=2, filter=wgts)
plot(soi)
lines(soif, lwd=2, col=4)
```



Fig. 2.11. The SOI series smoothed using (2.39) with $k=6$ (and half-weights at the ends). The insert shows the shape of the moving average ("boxcar") kernel [not drawn to scale] described in (2.41).


Fig. 2.12. Kernel smoother of the SOI. The insert shows the shape of the normal kernel [not drawn to scale].

Although the moving average smoother does a good job in highlighting the El Niño effect, it might be considered too choppy. We can obtain a smoother fit using the normal distribution for the weights, instead of boxcar-type weights of (2.39).

## Example 2.11 Kernel Smoothing

Kernel smoothing is a moving average smoother that uses a weight function, or kernel, to average the observations. Figure 2.12 shows kernel smoothing of the mortality series, where $m_{t}$ is now

$$
\begin{equation*}
m_{t}=\sum_{i=1}^{n} w_{i}(t) x_{i} \tag{2.40}
\end{equation*}
$$

where

$$
\begin{equation*}
w_{i}(t)=K\left(\frac{t-i}{b}\right) / \sum_{j=1}^{n} K\left(\frac{t-j}{b}\right) \tag{2.41}
\end{equation*}
$$

are the weights and $K(\cdot)$ is a kernel function. This estimator, which was originally explored by Parzen (1962) and Rosenblatt (1956b), is often called the Nadaraya-


Fig. 2.13. Locally weighted scatterplot smoothers (lowess) of the SOI series.

Watson estimator (Watson, 1966). In this example, and typically, the normal kernel, $K(z)=\frac{1}{\sqrt{2 \pi}} \exp \left(-z^{2} / 2\right)$, is used.

To implement this in R, use the ksmooth function where a bandwidth can be chosen. The wider the bandwidth, $b$, the smoother the result. From the R ksmooth help file: The kernels are scaled so that their quartiles (viewed as probability densities) are at $\pm 0.25 *$ bandwidth. For the standard normal distribution, the quartiles are $\pm .674$.

In our case, we are smoothing over time, which is of the form $t / 12$ for the SOI time series. In Figure 2.12, we used the value of $b=1$ to correspond to approximately smoothing over one year (recall that SOI has freq $=12$ so that the time scale is $t / 12$ ). Figure 2.12 can be reproduced in R as follows. plot(soi)
lines(ksmooth(time(soi), soi, "normal", bandwidth=1), lwd=2, col=4)

## Example 2.12 Lowess

Another approach to smoothing a time plot is nearest neighbor regression. The technique is based on $k$-nearest neighbors regression, wherein one uses only the data $\left\{x_{t-k / 2}, \ldots, x_{t}, \ldots, x_{t+k / 2}\right\}$ to predict $x_{t}$ via regression, and then sets $m_{t}=\hat{x}_{t}$.

Lowess is a method of smoothing that is rather complex, but the basic idea is close to nearest neighbor regression. Figure 2.13 shows smoothing of mortality using the R function lowess (see Cleveland, 1979). First, a certain proportion of nearest neighbors to $x_{t}$ are included in a weighting scheme; values closer to $x_{t}$ in time get more weight. Then, a robust weighted regression is used to predict $x_{t}$ and obtain the smoothed values $m_{t}$. The larger the fraction of nearest neighbors included, the smoother the fit will be. In Figure 2.13, one smoother uses 5\% of the data to obtain an estimate of the El Niño cycle of the data. In addition, a (negative) trend in SOI would indicate the long-term warming of the Pacific Ocean. To investigate this, we used a lowess with the default smoother span of $f=2 / 3$ of the data.

Figure 2.13 can be reproduced in R or S-PLUS as follows. plot(soi)
lines(lowess(soi, f=.05), lwd=2, col=4) \# El Nino cycle
lines(lowess(soi), lty=2, lwd=2, col=2) \# trend (using default span)


Fig. 2.14. Smooth of mortality as a function of temperature using lowess.

## Example 2.13 Smoothing One Series as a Function of Another

In addition to smoothing time plots, smoothing techniques can be applied to smoothing a time series as a function of another time series. We already used this idea in Example 2.7 to visualize the nonlinearity between Recruitment and SOI at various lags via lowess. In this example, we smooth the scatterplot of two contemporaneously measured time series, mortality as a function of temperature. In Example 2.2, we discovered a nonlinear relationship between mortality and temperature. Continuing along these lines, Figure 2.14 shows scatterplots of mortality, $M_{t}$, and temperature, $T_{t}$, along with $M_{t}$ smoothed as a function of $T_{t}$ using lowess Note that mortality increases at extreme temperatures, but in an asymmetric way; mortality is higher at colder temperatures than at hotter temperatures. The minimum mortality rate seems to occur at approximately $80^{\circ} \mathrm{F}$.

Figure 2.14 can be reproduced in R as follows using the defaults.
plot(tempr, cmort, xlab="Temperature", ylab="Mortality") lines(lowess(tempr, cmort))

## Problems

2.1 (Structural Model) For the Johnson \& Johnson data, say $y_{t}$, shown in Figure 1.1, let $x_{t}=\log \left(y_{t}\right)$. In this problem, we are going to fit a special type of structural model, $x_{t}=T_{t}+S_{t}+N_{t}$ where $T_{t}$ is a trend component, $S_{t}$ is a seasonal component, and $N_{t}$ is noise. In our case, time $t$ is in quarters $(1960.00,1960.25, \ldots)$ so one unit of time is a year.
(a) Fit the regression model

$$
x_{t}=\underbrace{\beta t}_{\text {trend }}+\underbrace{\alpha_{1} Q_{1}(t)+\alpha_{2} Q_{2}(t)+\alpha_{3} Q_{3}(t)+\alpha_{4} Q_{4}(t)}_{\text {seasonal }}+\underbrace{w_{t}}_{\text {noise }}
$$

where $Q_{i}(t)=1$ if time $t$ corresponds to quarter $i=1,2,3,4$, and zero otherwise. The $Q_{i}(t)$ 's are called indicator variables. We will assume for now that $w_{t}$ is a

Gaussian white noise sequence. Hint: Detailed code is given in Appendix R on page 186.
(b) If the model is correct, what is the estimated average annual increase in the logged earnings per share?
(c) If the model is correct, does the average logged earnings rate increase or decrease from the third quarter to the fourth quarter? And, by what percentage does it increase or decrease?
(d) What happens if you include an intercept term in the model in (a)? Explain why there was a problem.
(e) Graph the data, $x_{t}$, and superimpose the fitted values, say $\widehat{x}_{t}$, on the graph. Examine the residuals, $x_{t}-\widehat{x}_{t}$, and state your conclusions. Does it appear that the model fits the data well (do the residuals look white)?
2.2 For the mortality data examined in Example 2.2:
(a) Add another component to the regression in (2.21) that accounts for the particulate count four weeks prior; that is, add $P_{t-4}$ to the regression in (2.21). State your conclusion.
(b) Using AIC and BIC, is the model in (a) an improvement over the final model in Example 2.2?
2.3 Read the entire problem before you start because you can do parts (a) and (b) at the same time.
(a) Generate four series that are random walk with drift, (1.4), of length $n=100$ with $\delta=.01$ and $\sigma_{w}=1$. Call the data $x_{t}$ for $t=1, \ldots, 100$. Fit the regression $x_{t}=\beta t+w_{t}$ using least squares. Plot the data, the true mean function (i.e., $\left.\mu_{t}=.01 t\right)$ and the fitted line, $\widehat{x}_{t}=\widehat{\beta} t$, on the same graph.
(b) Generate four series of length $n=100$ that are linear trend plus noise, say $y_{t}=.01 t+w_{t}$, where $t$ and $w_{t}$ are as in part (a). Fit the regression $y_{t}=\beta t+w_{t}$ using least squares. Plot the data, the true mean function (i.e., $\mu_{t}=.01 t$ ) and the fitted line, $\widehat{y}_{t}=\widehat{\beta} t$, on the same graph.
(c) Comment (what did you learn from this assignment).

The following R code may be useful for doing (a) and (b) at the same time.

```
par(mfrow=c (4,2), mar=c(2.5,2.5,0,0)+.5, mgp=c(1.6,.6,0)) # set up
for (i in 1:4){
    x = ts(cumsum(rnorm(100,.01,1))) # data
    y = ts(.01*1:100 + rnorm(100))
    regx = lm(x~0+time(x), na.action=NULL) # regressions
    regy = lm(y~0+time(y), na.action=NULL)
    plot(x) # plots
    lines(.01*time(x), col="red", lty="dashed") # true mean
        abline(regx, col="blue") # fitted line
    plot(y)
        lines(.01*time(y), col="red", lty="dashed")
        abline(regy, col="blue")
}
```

2.4 Consider a process consisting of a linear trend with an additive noise term consisting of independent random variables $w_{t}$ with zero means and variances $\sigma_{w}^{2}$, that is,

$$
x_{t}=\beta_{0}+\beta_{1} t+w_{t},
$$

where $\beta_{0}, \beta_{1}$ are fixed constants.
(a) Prove $x_{t}$ is nonstationary.
(b) Prove that the first difference series $\nabla x_{t}=x_{t}-x_{t-1}$ is stationary by finding its mean and autocovariance function.
(c) Repeat part (b) if $w_{t}$ is replaced by a general stationary process, say $y_{t}$, with mean function $\mu_{y}$ and autocovariance function $\gamma_{y}(h)$. [Hint: See (2.28).]
2.5 Show (2.27) is stationary.
2.6 The glacial varve record plotted in Figure 2.6 exhibits some nonstationarity that can be improved by transforming to logarithms and some additional nonstationarity that can be corrected by differencing the logarithms.
(a) Argue that the glacial varves series, say $x_{t}$, exhibits heteroscedasticity by computing the sample variance over the first half and the second half of the data. Argue that the transformation $y_{t}=\log x_{t}$ stabilizes the variance over the series. Plot the histograms of $x_{t}$ and $y_{t}$ to see whether the approximation to normality is improved by transforming the data.
(b) Plot the series $y_{t}$. Do any time intervals, of the order 100 years, exist where one can observe behavior comparable to that observed in the global temperature records in Figure 1.3?
(c) Examine the sample ACF of $y_{t}$ and comment.
(d) Compute the difference $u_{t}=y_{t}-y_{t-1}$, examine its time plot and sample ACF, and argue that differencing the logged varve data produces a reasonably stationary series. Can you think of a practical interpretation for $u_{t}$ ? Hint: For $|p|$ close to zero, $\log (1+p) \approx p$; let $p=\left(y_{t}-y_{t-1}\right) / y_{t-1}$.
2.7 Use the three different smoothing techniques described in Example 2.10, Example 2.11, and Example 2.12, to estimate the trend in the global temperature series displayed in Figure 1.3. Comment.

## Chapter 3

## ARIMA Models

### 3.1 Introduction

In Chapters 1 and 2, we introduced autocorrelation and cross-correlation functions (ACFs and CCFs) as tools for clarifying relations that may occur within and between time series at various lags. In addition, we explained how to build linear models based on classical regression theory for exploiting the associations indicated by large values of the ACF or CCF.

Classical regression is often insufficient for explaining all of the interesting dynamics of a time series. For example, the ACF of the residuals of the simple linear regression fit to the global temperature data (see Example 2.4 of Chapter 2) reveals additional structure in the data that the regression did not capture. Instead, the introduction of correlation as a phenomenon that may be generated through lagged linear relations leads to proposing the autoregressive (AR) and moving average (MA) models. Often, these models are combined to form the autoregressive moving average (ARMA) model. Adding nonstationary models to the mix leads to the autoregressive integrated moving average (ARIMA) model popularized in the landmark work by Box and Jenkins (1970). Seasonal data, such as the data discussed in Example 1.1 and Example 1.4 lead to seasonal autoregressive integrated moving average (SARIMA) models. The Box-Jenkins method for identifying a plausible models is given in this chapter along with techniques for parameter estimation and forecasting for these models.

### 3.2 Autoregressive Moving Average Models

Autoregressive models are based on the idea that the current value of the series, $x_{t}$, can be explained as a function of $p$ past values, $x_{t-1}, x_{t-2}, \ldots, x_{t-p}$, where $p$ determines the number of steps into the past needed to forecast the current value. As a typical case, recall Example 1.8 in which data were generated using the model

$$
x_{t}=x_{t-1}-.90 x_{t-2}+w_{t}
$$

where $w_{t}$ is white Gaussian noise with $\sigma_{w}^{2}=1$. We have now assumed the current value is a particular linear function of past values. The regularity that persists in Figure 1.8 gives an indication that forecasting for such a model might be a distinct possibility, say, through some version such as

$$
x_{n+1}^{n}=x_{n}-.90 x_{n-1},
$$

where the quantity on the left-hand side denotes the forecast at the next period $n+1$ based on the observed data, $x_{1}, x_{2}, \ldots, x_{n}$. For example, the lagged scatterplot matrix for the Southern Oscillation Index (SOI), shown in Figure 2.7, gives a distinct indication that the values at lags 1,2 , and 12 , are linearly associated with the current value. We will make this notion more precise in our discussion of forecasting.

Definition 3.1 An autoregressive model of order $p$, abbreviated $\mathbf{A R}(p)$, is of the form

$$
\begin{equation*}
x_{t}=\phi_{1} x_{t-1}+\phi_{2} x_{t-2}+\cdots+\phi_{p} x_{t-p}+w_{t}, \tag{3.1}
\end{equation*}
$$

where $x_{t}$ is stationary, and $\phi_{1}, \phi_{2}, \ldots, \phi_{p}$ are constants $\left(\phi_{p} \neq 0\right)$. Although it is not necessary yet, we assume that $w_{t}$ is a Gaussian white noise series with mean zero and variance $\sigma_{w}^{2}$, unless otherwise stated. The mean of $x_{t}$ in (3.1) is zero. If the mean, $\mu$, of $x_{t}$ is not zero, replace $x_{t}$ by $x_{t}-\mu$ in (3.1),

$$
x_{t}-\mu=\phi_{1}\left(x_{t-1}-\mu\right)+\phi_{2}\left(x_{t-2}-\mu\right)+\cdots+\phi_{p}\left(x_{t-p}-\mu\right)+w_{t},
$$

or write

$$
\begin{equation*}
x_{t}=\alpha+\phi_{1} x_{t-1}+\phi_{2} x_{t-2}+\cdots+\phi_{p} x_{t-p}+w_{t}, \tag{3.2}
\end{equation*}
$$

where $\alpha=\mu\left(1-\phi_{1}-\cdots-\phi_{p}\right)$.
We note that (3.2) is similar to the regression model of $\S 2.2$, and hence the term auto (or self) regression. Some technical difficulties, however, develop from applying that model because the regressors, $x_{t-1}, \ldots, x_{t-p}$, are random components, whereas $z_{t}$ was assumed to be fixed. A useful form follows by using the backshift operator (2.30) to write the $\operatorname{AR}(p)$ model, (3.1), as

$$
\begin{equation*}
\left(1-\phi_{1} B-\phi_{2} B^{2}-\cdots-\phi_{p} B^{p}\right) x_{t}=w_{t}, \tag{3.3}
\end{equation*}
$$

or even more concisely as

$$
\begin{equation*}
\phi(B) x_{t}=w_{t} . \tag{3.4}
\end{equation*}
$$

## Example 3.1 The AR(1) Model

Consider the first-order model, $\operatorname{AR}(1)$, given by $x_{t}=\phi x_{t-1}+w_{t}$. Provided that $|\phi|<1$ we can represent an $\operatorname{AR}(1)$ model as a linear process given by ${ }^{1}$

$$
\begin{equation*}
x_{t}=\sum_{j=0}^{\infty} \phi^{j} w_{t-j} \tag{3.5}
\end{equation*}
$$

Representation (3.5) is called the stationary solution of the model. In fact, by simple substitution,

$$
\underbrace{\sum_{j=0}^{\infty} \phi^{j} w_{t-j}}_{x_{t}}=\phi(\underbrace{\sum_{k=0}^{\infty} \phi^{k} w_{t-1-k}}_{x_{t-1}})+w_{t} .
$$

Using (3.5), it is easy to see that the $\operatorname{AR}(1)$ process is stationary with mean

$$
E\left(x_{t}\right)=\sum_{j=0}^{\infty} \phi^{j} E\left(w_{t-j}\right)=0,
$$

and autocovariance function ( $h \geq 0$ ),

$$
\begin{align*}
& \gamma(h)=\operatorname{cov}\left(x_{t+h}, x_{t}\right)=E\left[\left(\sum_{j=0}^{\infty} \phi^{j} w_{t+h-j}\right)\left(\sum_{k=0}^{\infty} \phi^{k} w_{t-k}\right)\right] \\
& =E\left[\left(w_{t+h}+\cdots+\phi^{h} w_{t}+\phi^{h+1} w_{t-1}+\cdots\right)\left(w_{t}+\phi w_{t-1}+\cdots\right)\right]  \tag{3.6}\\
& =\sigma_{w}^{2} \sum_{j=0}^{\infty} \phi^{h+j} \phi^{j}=\sigma_{w}^{2} \phi^{h} \sum_{j=0}^{\infty} \phi^{2 j}=\frac{\sigma_{w}^{2} \phi^{h}}{1-\phi^{2}}, \quad h \geq 0 .
\end{align*}
$$

Recall that $\gamma(h)=\gamma(-h)$, so we will only exhibit the autocovariance function for $h \geq 0$. From (3.6), the ACF of an $\operatorname{AR}(1)$ is

$$
\begin{equation*}
\rho(h)=\frac{\gamma(h)}{\gamma(0)}=\phi^{h}, \quad h \geq 0 . \tag{3.7}
\end{equation*}
$$

## Example 3.2 The Sample Path of an AR(1) Process

Figure 3.1 shows a time plot of two $\operatorname{AR}(1)$ processes, one with $\phi=.9$ and one with $\phi=-.9$; in both cases, $\sigma_{w}^{2}=1$. In the first case, $\rho(h)=.9^{h}$, for $h \geq 0$, so observations close together in time are positively correlated with each other. This result means that observations at contiguous time points will tend to be close in value to each other; this fact shows up in the top of Figure 3.1 as a very smooth sample path for $x_{t}$. Now, contrast this with the case in which $\phi=-.9$, so that

[^7]

Fig. 3.1. Simulated $\operatorname{AR}(1)$ models: $\phi=.9$ (top); $\phi=-.9$ (bottom).
$\rho(h)=(-.9)^{h}$, for $h \geq 0$. This result means that observations at contiguous time points are negatively correlated but observations two time points apart are positively correlated. This fact shows up in the bottom of Figure 3.1, where, for example, if an observation, $x_{t}$, is positive, the next observation, $x_{t+1}$, is typically negative, and the next observation, $x_{t+2}$, is typically positive. Thus, in this case, the sample path is very choppy.

The following R code can be used to obtain a figure similar to Figure 3.1:

$$
\operatorname{par}(\operatorname{mfrow}=c(2,1))
$$

plot(arima.sim(list(order=c(1,0,0), ar=.9), n=100), ylab="x",
main=(expression(AR (1)~~~phi==+.9)))
plot(arima.sim(list(order=c(1,0,0), ar=-.9), n=100), ylab="x",
main=(expression(AR(1)~~~phi==-.9)))

As an alternative to the autoregressive representation in which the $x_{t}$ on the lefthand side of the equation are assumed to be combined linearly, the moving average model of order $q$, abbreviated as $\operatorname{MA}(q)$, assumes the white noise $w_{t}$ on the right-hand side of the defining equation are combined linearly to form the observed data.

Definition 3.2 The moving average model of order $q$, or MA(q) model, is defined to be

$$
\begin{equation*}
x_{t}=w_{t}+\theta_{1} w_{t-1}+\theta_{2} w_{t-2}+\cdots+\theta_{q} w_{t-q} \tag{3.8}
\end{equation*}
$$

where there are $q$ lags in the moving average and $\theta_{1}, \theta_{2}, \ldots, \theta_{q}\left(\theta_{q} \neq 0\right)$ are parameters. ${ }^{2}$ Although it is not necessary yet, we assume that $w_{t}$ is a Gaussian white noise series with mean zero and variance $\sigma_{w}^{2}$, unless otherwise stated.

[^8]

Fig. 3.2. Simulated MA(1) models: $\theta=.5$ (top); $\theta=-.5$ (bottom).

As in the $\operatorname{AR}(p)$ case, the $\operatorname{MA}(q)$ model may be written as

$$
\begin{equation*}
x_{t}=\left(1+\theta_{1} B+\theta_{2} B^{2}+\cdots+\theta_{q} B^{q}\right) w_{t}, \tag{3.9}
\end{equation*}
$$

or more concisely as

$$
\begin{equation*}
x_{t}=\theta(B) w_{t} \tag{3.10}
\end{equation*}
$$

Unlike the autoregressive process, the moving average process is stationary for any values of the parameters $\theta_{1}, \ldots, \theta_{q}$.

## Example 3.3 The MA(1) Process

Consider the MA(1) model $x_{t}=w_{t}+\theta w_{t-1}$. Then, $E\left(x_{t}\right)=0$,

$$
\gamma(h)= \begin{cases}\left(1+\theta^{2}\right) \sigma_{w}^{2} & h=0 \\ \theta \sigma_{w}^{2} & h=1 \\ 0 & h>1\end{cases}
$$

and the ACF is

$$
\rho(h)= \begin{cases}\frac{\theta}{\left(1+\theta^{2}\right)} & h=1, \\ 0 & h>1 .\end{cases}
$$

Note $|\rho(1)| \leq 1 / 2$ for all values of $\theta$ (Problem 3.1). Also, $x_{t}$ is correlated with $x_{t-1}$, but not with $x_{t-2}, x_{t-3}, \ldots$. Contrast this with the case of the $\operatorname{AR}(1)$ model in which the correlation between $x_{t}$ and $x_{t-k}$ is never zero. When $\theta=.5$, for example, $x_{t}$ and $x_{t-1}$ are positively correlated, and $\rho(1)=.4$. When $\theta=-.5, x_{t}$ and $x_{t-1}$ are negatively correlated, $\rho(1)=-4$. Figure 3.2 shows a time plot of these two
processes with $\sigma_{w}^{2}=1$. The series where $\theta=.5$ is smoother than the series where $\theta=-.5$.

A figure similar to Figure 3.2 can be created in R as follows:
$\operatorname{par}(m f r o w=c(2,1))$
plot(arima.sim(list(order $=c(\theta, \theta, 1), m a=.5), n=10 \theta), y l a b=" x "$,
main=(expression(MA(1)~~~theta==+.5)))
plot(arima.sim(list(order=c $(\theta, 0,1)$, ma=-.5), $n=10 \theta)$, ylab="x",
main=(expression(MA(1)~~~theta==-.5)))
We now proceed with the general development of mixed autoregressive moving average (ARMA) models for stationary time series.

Definition 3.3 A time series $\left\{x_{t} ; t=0, \pm 1, \pm 2, \ldots\right\}$ is ARMA $(p, q)$ if it is stationary and

$$
\begin{equation*}
x_{t}=\phi_{1} x_{t-1}+\cdots+\phi_{p} x_{t-p}+w_{t}+\theta_{1} w_{t-1}+\cdots+\theta_{q} w_{t-q} \tag{3.11}
\end{equation*}
$$

with $\phi_{p} \neq 0, \theta_{q} \neq 0$, and $\sigma_{w}^{2}>0$. The parameters $p$ and $q$ are called the autoregressive and the moving average orders, respectively. If $x_{t}$ has a nonzero mean $\mu$, we set $\alpha=\mu\left(1-\phi_{1}-\cdots-\phi_{p}\right)$ and write the model as

$$
\begin{equation*}
x_{t}=\alpha+\phi_{1} x_{t-1}+\cdots+\phi_{p} x_{t-p}+w_{t}+\theta_{1} w_{t-1}+\cdots+\theta_{q} w_{t-q} \tag{3.12}
\end{equation*}
$$

Although it is not necessary yet, we assume that $w_{t}$ is a Gaussian white noise series with mean zero and variance $\sigma_{w}^{2}$, unless otherwise stated.

The ARMA model may be seen as a regression of the present outcome $\left(x_{t}\right)$ on the past outcomes $\left(x_{t-1}, \ldots, x_{t-p}\right)$, with correlated errors. That is,

$$
x_{t}=\beta_{0}+\beta_{1} x_{t-1}+\cdots+\beta_{p} x_{t-p}+\epsilon_{t}
$$

where $\epsilon_{t}=w_{t}+\theta_{1} w_{t-1}+\cdots+\theta_{q} w_{t-q}$, although we call the regression parameters $\phi$ instead of $\beta$.

As previously noted, when $q=0$, the model is called an autoregressive model of order $p, \operatorname{AR}(p)$, and when $p=0$, the model is called a moving average model of order $q$, MA $(q)$. Using (3.3) and (3.9), the $\operatorname{ARMA}(p, q)$ model in (3.11) may be written in concise form as

$$
\begin{equation*}
\phi(B) x_{t}=\theta(B) w_{t} . \tag{3.13}
\end{equation*}
$$

The concise form of an ARMA model points to a potential problem in that we can unnecessarily complicate the model by multiplying both sides by another operator, say

$$
\eta(B) \phi(B) x_{t}=\eta(B) \theta(B) w_{t},
$$

without changing the dynamics. Consider the following example.

## Example 3.4 Parameter Redundancy

Consider a white noise process $x_{t}=w_{t}$. Equivalently, we can write this as $.5 x_{t-1}=$ $.5 w_{t-1}$ by shifting back one unit of time and multiplying by .5 . Now, subtract the two representations to obtain

$$
x_{t}-.5 x_{t-1}=w_{t}-.5 w_{t-1}
$$

or

$$
\begin{equation*}
x_{t}=.5 x_{t-1}-.5 w_{t-1}+w_{t}, \tag{3.14}
\end{equation*}
$$

which looks like an $\operatorname{ARMA}(1,1)$ model. Of course, $x_{t}$ is still white noise; nothing has changed in this regard [i.e., $x_{t}=w_{t}$ is the solution to (3.14)], but we have hidden the fact that $x_{t}$ is white noise because of the parameter redundancy or over-parameterization. Write the parameter redundant model in operator form as $\phi(B) x_{t}=\theta(B) w_{t}$, or

$$
(1-.5 B) x_{t}=(1-.5 B) w_{t} .
$$

Apply the operator $\phi(B)^{-1}=(1-.5 B)^{-1}$ to both sides to obtain

$$
x_{t}=(1-.5 B)^{-1}(1-.5 B) x_{t}=(1-.5 B)^{-1}(1-.5 B) w_{t}=w_{t},
$$

which is the original model.
Example 3.4 points out the need to be careful when fitting ARMA models to data. For example, if a process is truly white noise, it is possible to fit a significant $\operatorname{ARMA}(k, k)$ model to the data. That is, it is possible to obtain a seemingly complicated dynamic description of simple white noise. Consider the following example.

## Example 3.5 Parameter Redundancy (cont)

Although we have not yet discussed estimation, we present the following example demonstrating the problem of over parameterization. We generated 150 iid standard normals and then fit an $\operatorname{ARMA}(1,1)$ to the data. Note that $\widehat{\phi}=-.96$ and $\widehat{\theta}=.95$, and both are significant. Below is the R code (note that the estimate called 'intercept' is really the estimate of the mean).

```
set.seed(8675309) # Jenny, I got your number
x = rnorm(150) # generate iid N(0,1)s
arima(x, order=c(1,0,1)) # estimation
    Coefficients:
                        ar1 ma1 intercept <- misnomer
        -0.9595
```

    Thus, forgetting the (non-significant) mean estimate, the fitted model looks like
    $$
(1+.96 B) x_{t}=(1+.95 B) w_{t}
$$

## Definition 3.4 Causality and Invertibility

The causality and invertibility conditions are conditions on the parameters that ensure the present will not depend on the future. These conditions are obviously crucial for forecasting.

The causal form of the model is given by

$$
\begin{equation*}
x_{t}=\phi(B)^{-1} \theta(B) w_{t}=\psi(B) w_{t}=\sum_{j=0}^{\infty} \psi_{j} w_{t-j}, \tag{3.15}
\end{equation*}
$$

where $\psi(B)=\sum_{j=0}^{\infty} \psi_{j} B^{j}\left(\psi_{0}=1\right)$. Note that the parameters $\psi_{j}$ may be obtained by matching coefficients of $B$ in $\phi(B) \psi(B)=\theta(B)$.

The invertible form of the model is given by

$$
\begin{equation*}
w_{t}=\theta(B)^{-1} \phi(B) x_{t}=\pi(B) x_{t}=\sum_{j=0}^{\infty} \pi_{j} x_{t-j} . \tag{3.16}
\end{equation*}
$$

where $\pi(B)=\sum_{j=0}^{\infty} \pi_{j} B^{j}\left(\pi_{0}=1\right)$, assuming the representations are well-defined. Likewise, the parameters $\pi_{j}$ may be obtained by matching coefficients of $B$ in $\phi(B)=$ $\pi(B) \theta(B)$.

We note that it is not always possible to solve these relationships and some restrictions apply, as follows.

## Property 3.1 Causality and Invertibility (existence)

Let

$$
\phi(z)=1-\phi_{1} z-\cdots-\phi_{p} z^{p} \quad \text { and } \quad \theta(z)=1+\theta_{1} z+\cdots+\theta_{q} z^{q}
$$

be the AR and MA polynomials obtained by replacing the backshift operator $B$ in (3.3) and (3.9) by a complex number $z$.
$\operatorname{An} \operatorname{ARMA}(p, q)$ model is causal if and only if $\phi(z) \neq 0$ for $|z| \leq 1$. The coefficients of the linear process given in (3.15) can be determined by solving $\left(\psi_{0}=1\right)$

$$
\psi(z)=\sum_{j=0}^{\infty} \psi_{j} z^{j}=\frac{\theta(z)}{\phi(z)}, \quad|z| \leq 1 .^{*}
$$

An $\operatorname{ARMA}(p, q)$ model is invertible if and only if $\theta(z) \neq 0$ for $|z| \leq 1$. The coefficients $\pi_{j}$ of $\pi(B)$ given in (3.16) can be determined by solving ( $\pi_{0}=1$ )

$$
\pi(z)=\sum_{j=0}^{\infty} \pi_{j} z^{j}=\frac{\phi(z)}{\theta(z)}, \quad|z| \leq 1 .^{\dagger}
$$

We demonstrate the property in the following example.

## Example 3.6 Parameter Redundancy, Causality, Invertibility

Consider the process

$$
x_{t}=.4 x_{t-1}+.45 x_{t-2}+w_{t}+w_{t-1}+.25 w_{t-2},
$$

or, in operator form,

$$
\left(1-.4 B-.45 B^{2}\right) x_{t}=\left(1+B+.25 B^{2}\right) w_{t} .
$$

At first, $x_{t}$ appears to be an $\operatorname{ARMA}(2,2)$ process. But notice that

[^9]$$
\phi(B)=1-.4 B-.45 B^{2}=(1+.5 B)(1-.9 B)
$$
and
$$
\theta(B)=\left(1+B+.25 B^{2}\right)=(1+.5 B)^{2}
$$
have a common factor that can be canceled. After cancellation, the operators are $\phi(B)=(1-.9 B)$ and $\theta(B)=(1+.5 B)$, so the model is an ARMA $(1,1)$ model, $(1-.9 B) x_{t}=(1+.5 B) w_{t}$, or
\[

$$
\begin{equation*}
x_{t}=.9 x_{t-1}+.5 w_{t-1}+w_{t} . \tag{3.17}
\end{equation*}
$$

\]

The model is causal because $\phi(z)=(1-.9 z)=0$ when $z=10 / 9$, which is outside the unit circle. The model is also invertible because the root of $\theta(z)=$ $(1+.5 z)$ is $z=-2$, which is outside the unit circle.

To write the model as a linear process, we can obtain the $\psi$-weights using Property 3.1, $\phi(z) \psi(z)=\theta(z)$, or

$$
(1-.9 z)\left(1+\psi_{1} z+\psi_{2} z^{2}+\cdots+\psi_{j} z^{j}+\cdots\right)=1+.5 z .
$$

Rearranging, we get

$$
1+\left(\psi_{1}-.9\right) z+\left(\psi_{2}-.9 \psi_{1}\right) z^{2}+\cdots+\left(\psi_{j}-.9 \psi_{j-1}\right) z^{j}+\cdots=1+.5 z
$$

The coefficients of $z$ on the left and right sides must be the same, so we get $\psi_{1}-.9=.5$ or $\psi_{1}=1.4$, and $\psi_{j}-.9 \psi_{j-1}=0$ for $j>1$. Thus, $\psi_{j}=1.4(.9)^{j-1}$ for $j \geq 1$ and (3.17) can be written as

$$
x_{t}=w_{t}+1.4 \sum_{j=1}^{\infty} .9^{j-1} w_{t-j}
$$

The values of $\psi_{j}$ may be calculated in R as follows:

```
ARMAtoMA(ar = .9, ma = .5, 10) # first 10 psi-weights
```

$\begin{array}{llllllllllllllllllllll}{[1.40} & 1.26 & 1.13 & 1.02 & 0.92 & 0.83 & 0.74 & 0.67 & 0.60 & 0.54\end{array}$

The invertible representation using Property 3.1 is obtained by matching coefficients in $\theta(z) \pi(z)=\phi(z)$,

$$
(1+.5 z)\left(1+\pi_{1} z+\pi_{2} z^{2}+\pi_{3} z^{3}+\cdots\right)=1-.9 z
$$

In this case, the $\pi$-weights are given by $\pi_{j}=(-1)^{j} 1.4(.5)^{j-1}$, for $j \geq 1$, and hence, we can also write (3.17) as

$$
x_{t}=-1.4 \sum_{j=1}^{\infty}(-.5)^{j-1} x_{t-j}+w_{t} .
$$

The values of $\pi_{j}$ may be calculated in R as follows by reversing the roles of $w_{t}$ and $x_{t}$; i.e., write the model as $w_{t}=-.5 w_{t-1}+x_{t}-.9 x_{t-1}$ :
ARMAtoMA (ar $=-.5$, ma $=-.9,10) \quad$ \# first 10 pi-weights
[1] -1.400 . $700-.350$. $175-.087$. $044-.022$. $011-.006$. 003


Fig. 3.3. Causal region for an $\operatorname{AR}(2)$ in terms of the parameters.

## Example 3.7 Causal Conditions for an AR(2) Process

For an $\operatorname{AR}(1)$ model, $(1-\phi B) x_{t}=w_{t}$, to be causal, we must have $\phi(z) \neq 0$ for $|z| \leq 1$. If we solve $\phi(z)=1-\phi z=0$, we find that the root (or zero) occurs at $z_{0}=1 / \phi$, so that $\left|z_{0}\right|>1$ only if $|\phi|<1$. For higher order models, the relationship between roots and parameters is not so simple.

For example, the $\operatorname{AR}(2)$ model, $\left(1-\phi_{1} B-\phi_{2} B^{2}\right) x_{t}=w_{t}$, is causal when the two roots of $\phi(z)=1-\phi_{1} z-\phi_{2} z^{2}$ lie outside of the unit circle. That is, if $z_{1}$ and $z_{2}$ are the roots, then $\left|z_{1}\right|>1$ and $\left|z_{2}\right|>1$. Using the quadratic formula, this requirement can be written as

$$
\left|\frac{\phi_{1} \pm \sqrt{\phi_{1}^{2}+4 \phi_{2}}}{-2 \phi_{2}}\right|>1
$$

The roots of $\phi(z)$ may be real and distinct, real and equal, or a complex conjugate pair. In terms of the coefficients, the equivalent condition is

$$
\begin{equation*}
\phi_{1}+\phi_{2}<1, \quad \phi_{2}-\phi_{1}<1, \quad \text { and } \quad\left|\phi_{2}\right|<1 . \tag{3.18}
\end{equation*}
$$

This causality condition specifies a triangular region in the parameter space; see Figure 3.3.

### 3.3 Autocorrelation and Partial Autocorrelation

We begin by exhibiting the ACF of an $\mathrm{MA}(q)$ process.

## Example 3.8 ACF of an MA(q)

The model is $x_{t}=\theta(B) w_{t}$, where $\theta(B)=1+\theta_{1} B+\cdots+\theta_{q} B^{q}$. Because $x_{t}$ is a finite linear combination of white noise terms, the process is stationary with mean

$$
E\left(x_{t}\right)=\sum_{j=0}^{q} \theta_{j} E\left(w_{t-j}\right)=0
$$

where we have written $\theta_{0}=1$, and with autocovariance function

$$
\begin{align*}
\gamma(h)=\operatorname{cov}\left(x_{t+h}, x_{t}\right) & =\operatorname{cov}\left(\sum_{j=0}^{q} \theta_{j} w_{t+h-j}, \sum_{k=0}^{q} \theta_{k} w_{t-k}\right) \\
& = \begin{cases}\sigma_{w}^{2} \sum_{j=0}^{q-h} \theta_{j} \theta_{j+h}, & 0 \leq h \leq q \\
0 & h>q .\end{cases} \tag{3.19}
\end{align*}
$$

Recall that $\gamma(h)=\gamma(-h)$, so we will only display the values for $h \geq 0$. The cutting off of $\gamma(h)$ after $q$ lags is the signature of the MA $(q)$ model. Dividing (3.19) by $\gamma(0)$ yields the ACF of an MA $(q)$ :

$$
\rho(h)= \begin{cases}\frac{\sum_{j=0}^{q-h} \theta_{j} \theta_{j+h}}{1+\theta_{1}^{2}+\cdots+\theta_{q}^{2}} & 1 \leq h \leq q  \tag{3.20}\\ 0 & h>q\end{cases}
$$

## Example 3.9 ACF of an $\operatorname{AR}(p)$ and $\operatorname{ARMA}(p, q)$

For an $\operatorname{AR}(p)$ or $\operatorname{ARMA}(p, q)$ model, $\phi(B) x_{t}=\theta(B) w_{t}$, write it as

$$
x_{t}=\phi(B)^{-1} \theta(B) w_{t}=\psi(B) w_{t}
$$

or

$$
\begin{equation*}
x_{t}=\sum_{j=0}^{\infty} \psi_{j} w_{t-j} \tag{3.21}
\end{equation*}
$$

It follows immediately that $E\left(x_{t}\right)=0$. Also, the autocovariance function of $x_{t}$ can be written as

$$
\begin{equation*}
\gamma(h)=\operatorname{cov}\left(x_{t+h}, x_{t}\right)=\sigma_{w}^{2} \sum_{j=0}^{\infty} \psi_{j} \psi_{j+h}, \quad h \geq 0 \tag{3.22}
\end{equation*}
$$

so that the ACF is given by

$$
\begin{equation*}
\rho(h)=\frac{\sum_{j=0}^{\infty} \psi_{j} \psi_{j+h}}{\sum_{j=0}^{\infty} \psi_{j}^{2}}, \quad h \geq 0 \tag{3.23}
\end{equation*}
$$

Unlike the $\operatorname{MA}(q)$, the $\operatorname{ACF}$ of an $\operatorname{AR}(p)$ or an $\operatorname{ARMA}(p, q)$ does not cut off at any lag, so using the ACF to help identify the order of an AR or ARMA is difficult. Also, (3.23) is not appealing in that it provides little information about the appearance of the ACF of various models.

## Example 3.10 The ACF of an AR(2) Process

Suppose $x_{t}=\phi_{1} x_{t-1}+\phi_{2} x_{t-2}+w_{t}$ is a causal AR(2) process. Multiply each side of the model by $x_{t-h}$ for $h>0$, and take expectation:

$$
E\left(x_{t} x_{t-h}\right)=\phi_{1} E\left(x_{t-1} x_{t-h}\right)+\phi_{2} E\left(x_{t-2} x_{t-h}\right)+E\left(w_{t} x_{t-h}\right) .
$$

The result is

$$
\begin{equation*}
\gamma(h)=\phi_{1} \gamma(h-1)+\phi_{2} \gamma(h-2), \quad h=1,2, \ldots . \tag{3.24}
\end{equation*}
$$

In (3.24), we used the fact that $E\left(x_{t}\right)=0$ and for $h>0, E\left(w_{t} x_{t-h}\right)=0$ because, by causality, $x_{t-h}$ does not depend on future errors. Divide (3.24) through by $\gamma(0)$ to obtain a recursion for the ACF:

$$
\begin{equation*}
\rho(h)-\phi_{1} \rho(h-1)-\phi_{2} \rho(h-2)=0, \quad h=1,2, \ldots . \tag{3.25}
\end{equation*}
$$

The initial conditions are $\rho(0)=1$ and $\rho(-1)=\phi_{1} /\left(1-\phi_{2}\right)$, which is obtained by evaluating (3.25) for $h=1$ and noting that $\rho(1)=\rho(-1)$.

Equations such as (3.25) are called difference equations, and the solutions are fairly simple expressions. First, the polynomial associated with (3.25) is

$$
\phi(z)=1-\phi_{1} z-\phi_{2} z^{2}
$$

where the power of $z$ is the power of the backshift, $B$; i.e., (3.25) is $\left(1-\phi_{1} B-\right.$ $\left.\phi_{2} B^{2}\right) \rho(h)=0$. In general, $z$ is a complex number. Let $z_{1}$ and $z_{2}$ be the roots (or zeros) of the associated polynomial, i.e., $\phi\left(z_{1}\right)=\phi\left(z_{2}\right)=0$. For a causal model, the roots are outside the unit circle: $\left|z_{1}\right|>1$ and $\left|z_{2}\right|>1$. Now, consider the solutions:
(i) When $z_{1}$ and $z_{2}$ are distinct, then

$$
\rho(h)=c_{1} z_{1}^{-h}+c_{2} z_{2}^{-h}
$$

so $\rho(h) \rightarrow 0$ exponentially fast as $h \rightarrow \infty$. The constants $c_{1}$ and $c_{2}$ are obtained by solving for them using the initial conditions given above. For example, when $h=0$, we have $1=c_{1}+c_{2}$, and so on.
(ii) When $z_{1}=z_{2}\left(=z_{0}\right)$ are equal (and hence real), then

$$
\rho(h)=z_{0}^{-h}\left(c_{1}+c_{2} h\right),
$$

so $\rho(h) \rightarrow 0$ exponentially fast as $h \rightarrow \infty$.
In case (i) with complex roots, $z_{2}=\bar{z}_{1}$ are a complex conjugate pair, and $c_{2}=\bar{c}_{1}$ [because $\rho(h)$ is real], and

$$
\rho(h)=c_{1} z_{1}^{-h}+\bar{c}_{1} \bar{z}_{1}^{-h} .
$$

Write $c_{1}$ and $z_{1}$ in polar coordinates, for example, $z_{1}=\left|z_{1}\right| e^{i \theta}$, where $\theta$ is the angle whose tangent is the ratio of the imaginary part and the real part of $z_{1}$ (sometimes called $\arg \left(z_{1}\right)$; the range of $\theta$ is $\left.[-\pi, \pi]\right)$. Then, using the fact that $e^{i \alpha}+e^{-i \alpha}=2 \cos (\alpha)$, the solution has the form

$$
\rho(h)=a\left|z_{1}\right|^{-h} \cos (h \theta+b),
$$

where $a$ and $b$ are determined by the initial conditions. Again, $\rho(h)$ dampens to zero exponentially fast as $h \rightarrow \infty$, but it does so in a sinusoidal fashion. The implication of this result is shown in Example 3.11.

## Example 3.11 An AR(2) with Complex Roots

Figure 3.4 shows $n=144$ observations from the $\operatorname{AR}(2)$ model

$$
x_{t}=1.5 x_{t-1}-.75 x_{t-2}+w_{t},
$$

with $\sigma_{w}^{2}=1$, and with complex roots chosen so the process exhibits pseudocyclic behavior at the rate of one cycle every 12 time points. The autoregressive polynomial for this model is $\phi(z)=1-1.5 z+.75 z^{2}$. The roots of $\phi(z)$ are $1 \pm i / \sqrt{3}$, and $\theta=\tan ^{-1}(1 / \sqrt{3})=2 \pi / 12$ radians per unit time. To convert the angle to cycles per unit time, divide by $2 \pi$ to get $1 / 12$ cycles per unit time. The ACF for this model is shown in $\S 3.4$, Figure 3.5 .
To calculate the roots of the polynomial and solve for arg in R:

```
z = c(1,-1.5,.75) # coefficients of the polynomial
(a = polyroot(z)[1]) # print one root: 1+0.57735i = 1 + i/sqrt(3)
arg = Arg(a)/(2*pi) # arg in cycles/pt
1/arg # = 12, the pseudo period
```

To reproduce Figure 3.4:

```
set.seed(8675309)
ar2 = arima.sim(list(order=c(2,0,0), ar=c(1.5,-.75)), n = 144)
plot(ar2, axes=FALSE, xlab="Time")
axis(2); axis(1, at=seq(0,144,by=12)); box()
abline(v=seq(0,144,by=12), lty=2)
```

To calculate and display the ACF for this model:

```
ACF = ARMAacf(ar=c (1.5,-.75), ma=0, 50)
plot(ACF, type='h", xlab="lag")
abline(h=0)
```

In general, the behavior of the $\operatorname{ACF}$ of an $\operatorname{AR}(p)$ or an $\operatorname{ARMA}(p, q)$ when $p \geq 2$ will be similar to the $\operatorname{AR}(2)$ case. When $p=1$, the behavior is like the $\operatorname{AR}(1)$ case.

## Example 3.12 The ACF of an ARMA $(1,1)$

Consider the ARMA $(1,1)$ process $x_{t}=\phi x_{t-1}+\theta w_{t-1}+w_{t}$, where $|\phi|<1$. Using the theory of difference equations, we can show that the ACF is given by

$$
\begin{equation*}
\rho(h)=\frac{(1+\theta \phi)(\phi+\theta)}{1+2 \theta \phi+\theta^{2}} \phi^{h-1}, \quad h \geq 1 . \tag{3.26}
\end{equation*}
$$

Notice that the general pattern of $\rho(h)$ in (3.26) is not different from that of an AR(1) given in (3.7). Hence, it is unlikely that we will be able to tell the difference between an ARMA $(1,1)$ and an AR(1) based solely on an ACF estimated from a sample. This consideration will lead us to the partial autocorrelation function.


Fig. 3.4. Simulated $\operatorname{AR}(2)$ model, $n=144$ with $\phi_{1}=1.5$ and $\phi_{2}=-.75$.

## The Partial Autocorrelation Function (PACF)

In (3.20), we saw that for MA $(q)$ models, the ACF will be zero for lags greater than $q$. Moreover, because $\theta_{q} \neq 0$, the ACF will not be zero at lag $q$. Thus, the ACF provides a considerable amount of information about the order of the dependence when the process is a moving average process.

If the process, however, is ARMA or AR, the ACF alone tells us little about the orders of dependence. Hence, it is worthwhile pursuing a function that will behave like the ACF of MA models, but for AR models, namely, the partial autocorrelation function (PACF).

Recall that if $X, Y$, and $Z$ are random variables, then the partial correlation between $X$ and $Y$ given $Z$ is obtained by regressing $X$ on $Z$ to obtain the predictor $\widehat{X}$, regressing $Y$ on $Z$ to obtain $\widehat{Y}$, and then calculating

$$
\rho_{X Y \mid Z}=\operatorname{corr}\{X-\widehat{X}, Y-\widehat{Y}\} .
$$

The idea is that $\rho_{X Y \mid Z}$ measures the correlation between $X$ and $Y$ with the linear effect of $Z$ removed (or partialled out). If the variables are multivariate normal, then this definition coincides with $\rho_{X Y \mid Z}=\operatorname{corr}(X, Y \mid Z)$.

To motivate the idea of partial autocorrelation, consider a causal AR(1) model, $x_{t}=\phi x_{t-1}+w_{t}$. Then,

$$
\begin{aligned}
\gamma_{x}(2)=\operatorname{cov}\left(x_{t}, x_{t-2}\right) & =\operatorname{cov}\left(\phi x_{t-1}+w_{t}, x_{t-2}\right) \\
& =\operatorname{cov}\left(\phi^{2} x_{t-2}+\phi w_{t-1}+w_{t}, x_{t-2}\right)=\phi^{2} \gamma_{x}(0) .
\end{aligned}
$$

This result follows from causality because $x_{t-2}$ involves $\left\{w_{t-2}, w_{t-3}, \ldots\right\}$, which are all uncorrelated with $w_{t}$ and $w_{t-1}$. The correlation between $x_{t}$ and $x_{t-2}$ is not zero, as it would be for an MA(1), because $x_{t}$ is dependent on $x_{t-2}$ through $x_{t-1}$. Suppose we break this chain of dependence by removing (or partialling out) the effect of $x_{t-1}$. That is, we consider the correlation between $x_{t}-\phi x_{t-1}$ and $x_{t-2}-\phi x_{t-1}$, because it is the correlation between $x_{t}$ and $x_{t-2}$ with the linear dependence of each on $x_{t-1}$ removed. In this way, we have broken the dependence chain between $x_{t}$ and $x_{t-2}$. In fact,

$$
\operatorname{cov}\left(x_{t}-\phi x_{t-1}, x_{t-2}-\phi x_{t-1}\right)=\operatorname{cov}\left(w_{t}, x_{t-2}-\phi x_{t-1}\right)=0 .
$$

Hence, the tool we need is partial autocorrelation, which is the correlation between $x_{s}$ and $x_{t}$ with the linear effect of everything "in the middle" removed.

Definition 3.5 The partial autocorrelation function (PACF) of a stationary process, $x_{t}$, denoted $\phi_{h h}$, for $h=1,2, \ldots$, is

$$
\begin{equation*}
\phi_{11}=\operatorname{corr}\left(x_{1}, x_{0}\right)=\rho(1) \tag{3.27}
\end{equation*}
$$

and

$$
\begin{equation*}
\phi_{h h}=\operatorname{corr}\left(x_{h}-\widehat{x}_{h}, x_{0}-\widehat{x}_{0}\right), \quad h \geq 2, \tag{3.28}
\end{equation*}
$$

where $\widehat{x}_{h}$ is the regression of $x_{h}$ on $\left\{x_{1}, x_{2}, \ldots, x_{h-1}\right\}$ and $\widehat{x}_{0}$ is the regression of $x_{0}$ on $\left\{x_{1}, x_{2}, \ldots, x_{h-1}\right\}$.

Thus, due to the stationarity, the PACF, $\phi_{h h}$, is the correlation between $x_{t+h}$ and $x_{t}$ with the linear dependence of everything between them, namely $\left\{x_{t+1}, \ldots, x_{t+h-1}\right\}$, on each, removed.

## Example 3.13 The PACF of an $\operatorname{AR}(p)$

The model can be written as

$$
x_{t+h}=\sum_{j=1}^{p} \phi_{j} x_{t+h-j}+w_{t+h}
$$

where the roots of $\phi(z)$ are outside the unit circle. When $h>p$, the regression of $x_{t+h}$ on $\left\{x_{t+1}, \ldots, x_{t+h-1}\right\}$, is

$$
\widehat{x}_{t+h}=\sum_{j=1}^{p} \phi_{j} x_{t+h-j} .
$$

Although we have not proved this result, it should be obvious that it is so. Thus, when $h>p$,

$$
\phi_{h h}=\operatorname{corr}\left(x_{t+h}-\widehat{x}_{t+h}, x_{t}-\widehat{x}_{t}\right)=\operatorname{corr}\left(w_{t+h}, x_{t}-\widehat{x}_{t}\right)=0,
$$

because, by causality, $x_{t}-\widehat{x}_{t}$ depends only on $\left\{w_{t+h-1}, w_{t+h-2}, \ldots\right\}$. When $h \leq p$, $\phi_{p p}$ is not zero, and $\phi_{11}, \ldots, \phi_{p-1, p-1}$ are not necessarily zero. We will see later that, in fact, $\phi_{p p}=\phi_{p}$. Figure 3.5 shows the ACF and the PACF of the AR(2) model presented in Example 3.11.

To reproduce Figure 3.5 in R, use the following commands:

```
ACF = ARMAacf(ar=c (1.5,-.75), ma=0, 24)[-1]
PACF = ARMAacf(ar=c(1.5,-.75), ma=0, 24, pacf=TRUE)
par(mfrow=c(1,2))
plot(ACF, type="h", xlab="lag", ylim=c(-.8,1))
abline(h=0)
plot(PACF, type="h", xlab="lag", ylim=c(-.8,1))
abline(h=0)
```

We also have the following large sample result for the PACF, which may be compared to the similar result for the ACF given in Property 1.2.


Fig. 3.5. The ACF and PACF of an $\operatorname{AR}(2)$ model with $\phi_{1}=1.5$ and $\phi_{2}=-.75$.
Table 3.1. Behavior of the ACF and PACF for ARMA Models

|  | $\operatorname{AR}(p)$ | MA $(q)$ | ARMA $(p, q)$ |
| :---: | :---: | :---: | :---: |
| ACF | Tails off | Cuts off <br> after lag $q$ | Tails off |
| PACF | Cuts off <br> after lag $p$ | Tails off | Tails off |

## Property 3.2 Large Sample Distribution of the PACF

If the time series is a causal $A R(p)$ process and the sample size $n$ is large, then $\sqrt{n} \widehat{\phi}_{h h}$ is approximately $N(0,1)$, for $h>p$. This result also holds for $p=0$, wherein the process is white noise.

## Example 3.14 The PACF of an MA(q)

For an $\operatorname{MA}(q)$, we can write $x_{t}=-\sum_{j=1}^{\infty} \pi_{j} x_{t-j}+w_{t}$. Moreover, no finite representation exists. From this result, it should be apparent that the PACF will never cut off, as in the case of an $\operatorname{AR}(p)$. For an $\operatorname{MA}(1), x_{t}=w_{t}+\theta w_{t-1}$, with $|\theta|<1$, it can be shown that

$$
\phi_{h h}=-\frac{(-\theta)^{h}\left(1-\theta^{2}\right)}{1-\theta^{2(h+1)}}, \quad h \geq 1 .
$$

We do not have to compute the PACF by performing numerous regressions first. The computations are done via a recursive formula that we will discuss later in Property 3.3.

The PACF for MA models behaves much like the ACF for AR models. Also, the PACF for AR models behaves much like the ACF for MA models. Because an invertible ARMA model has an infinite AR representation, the PACF will not cut off. We may summarize these results in Table 3.1.


Fig. 3.6. ACF and PACF of the Recruitment series. Note that the lag axes are in terms of season (12 months in this case).

## Example 3.15 Preliminary Analysis of the Recruitment Series

We consider the problem of modeling the Recruitment series shown in Figure 1.5. There are 453 months of observed recruitment ranging over the years 1950-1987. The ACF and the PACF given in Figure 3.6 are consistent with the behavior of an $\operatorname{AR}(2)$. The ACF has cycles corresponding roughly to a 12 -month period, and the PACF has large values for $h=1,2$ and then is essentially zero for higher order lags. Based on Table 3.1, these results suggest that a second-order $(p=2)$ autoregressive model might provide a good fit. Although we will discuss estimation in detail in $\S 3.6$, we ran a regression (see $\S 2.2$ ) using the data triplets $\left\{\left(x ; z_{1}, z_{2}\right)\right.$ : $\left.\left(x_{3} ; x_{2}, x_{1}\right),\left(x_{4} ; x_{3}, x_{2}\right), \ldots,\left(x_{453} ; x_{452}, x_{451}\right)\right\}$ to fit the model

$$
x_{t}=\phi_{0}+\phi_{1} x_{t-1}+\phi_{2} x_{t-2}+w_{t}
$$

for $t=3,4, \ldots, 453$. The values of the estimates were $\widehat{\phi}_{0}=6.74_{(1.11)}$, $\widehat{\phi}_{1}=1.35_{(.04)}, \widehat{\phi}_{2}=-.46_{(.04)}$, and $\widehat{\sigma}_{w}^{2}=89.72$, where the estimated standard errors are in parentheses.

The following $R$ code can be used for this analysis. We use the script acf2 to print and plot the ACF and PACF; see Appendix R for details.

```
acf2(rec, 48) # will produce values and a graphic
(regr = ar.ols(rec, order=2, demean=FALSE, intercept=TRUE))
regr$asy.se.coef # standard errors of the estimates
```


### 3.4 Estimation

Throughout this section, we assume we have $n$ observations, $x_{1}, \ldots, x_{n}$, from a causal and invertible Gaussian $\operatorname{ARMA}(p, q)$ process in which, initially, the order parameters, $p$ and $q$, are known. Our goal is to estimate the parameters, $\phi_{1}, \ldots, \phi_{p}, \theta_{1}, \ldots, \theta_{q}$, and $\sigma_{w}^{2}$. We will discuss the problem of determining $p$ and $q$ later in this section.

We begin with method of moments estimators. The idea behind these estimators is that of equating population moments to sample moments and then solving for the parameters in terms of the sample moments. We immediately see that, if $E\left(x_{t}\right)=$ $\mu$, then the method of moments estimator of $\mu$ is the sample average, $\bar{x}$. Thus, while discussing method of moments, we will assume $\mu=0$. Although the method of moments can produce good estimators, they can sometimes lead to suboptimal estimators. We first consider the case in which the method leads to optimal (efficient) estimators, that is, $\operatorname{AR}(p)$ models.

When the process is $\operatorname{AR}(p)$,

$$
x_{t}=\phi_{1} x_{t-1}+\cdots+\phi_{p} x_{t-p}+w_{t}
$$

similar to Example 3.10, we have the following result:
Definition 3.6 The Yule-Walker equations are given by

$$
\begin{align*}
\rho(h) & =\phi_{1} \rho(h-1)+\cdots+\phi_{p} \rho(h-p), \quad h=1,2, \ldots, p,  \tag{3.29}\\
\sigma_{w}^{2} & =\gamma(0)\left[1-\phi_{1} \rho(1)-\cdots-\phi_{p} \rho(p)\right] . \tag{3.30}
\end{align*}
$$

The estimators obtained by replacing $\gamma(0)$ with its estimate, $\widehat{\gamma}(0)$ and $\rho(h)$ with its estimate, $\widehat{\rho}(h)$, are called the Yule-Walker estimators. $\operatorname{For} \operatorname{AR}(p)$ models, if the sample size is large, the Yule-Walker estimators are approximately normally distributed, and $\widehat{\sigma}_{w}^{2}$ is close to the true value of $\sigma_{w}^{2}$.

## Example 3.16 Yule-Walker Estimation for an AR(2) Process

The data shown in Figure 3.4 were $n=144$ simulated observations from the $\operatorname{AR}(2)$ model $x_{t}=1.5 x_{t-1}-.75 x_{t-2}+w_{t}$, where $w_{t} \sim \operatorname{iid} \mathrm{~N}(0,1)$. Using the same simulated data, we have
ar.yw(ar2, order=2)

```
    Coefficients:
            1 2
    1.4471-0.7160
    sigma^2 estimated as 1.561
```


## Example 3.17 Yule-Walker Estimation of the Recruitment Series

In Example 3.15 we fit an $\operatorname{AR}(2)$ model to the recruitment series using regression. Below are the results of fitting the same model using Yule-Walker estimation in R, which are nearly identical to the values in Example 3.15.

```
rec.yw = ar.yw(rec, order=2)
rec.yw$x.mean # = 62.26 (mean estimate)
rec.yw$ar # = 1.33, -. 44 (parameter estimates)
sqrt(diag(rec.yw$asy.var.coef)) # = .04, .04 (standard errors)
rec.yw$var.pred # = 94.80 (error variance estimate)
```

In the case of $\operatorname{AR}(p)$ models, the Yule-Walker estimators are optimal estimators, but this is not true for $\operatorname{MA}(q)$ or $\operatorname{ARMA}(p, q)$ models. $\operatorname{AR}(p)$ models are linear models, and the Yule-Walker estimators are essentially least squares estimators. If we use method of moments for MA or ARMA models, we will not get optimal estimators because such processes are nonlinear in the parameters.

## Example 3.18 Method of Moments Estimation for an MA(1)

Consider the MA(1) model, $x_{t}=w_{t}+\theta w_{t-1}$, where $|\theta|<1$. The model can then be written as

$$
x_{t}=-\sum_{j=1}^{\infty}(-\theta)^{j} x_{t-j}+w_{t}
$$

which is nonlinear in $\theta$. The first two population autocovariances are $\gamma(0)=$ $\sigma_{w}^{2}\left(1+\theta^{2}\right)$ and $\gamma(1)=\sigma_{w}^{2} \theta$, so the estimate of $\theta$ is found by solving:

$$
\widehat{\rho}(1)=\frac{\widehat{\gamma}(1)}{\widehat{\gamma}(0)}=\frac{\widehat{\theta}}{1+\widehat{\theta}^{2}}
$$

Two solutions exist, so we would pick the invertible one. If $|\widehat{\rho}(1)| \leq \frac{1}{2}$, the solutions are real, otherwise, a real solution does not exist. Even though $|\rho(1)|<\frac{1}{2}$ for an invertible MA(1), it may happen that $|\widehat{\rho}(1)| \geq \frac{1}{2}$ because it is an estimator. For example, the following simulation in R produces a value of $\widehat{\rho}(1)=.507$ when the true value is $\rho(1)=.9 /\left(1+.9^{2}\right)=.497$.
set.seed(2)
ma1 $=$ arima. $\operatorname{sim}($ list (order $=c(0, \theta, 1)$, ma $=0.9), n=50)$
$\operatorname{acf}(\mathrm{ma1}, \mathrm{plot}=\mathrm{FALSE})[1] \quad \#=.507$ ( lag 1 sample ACF)
The preferred method of estimation is maximum likelihood estimation (MLE), which determines the values of the parameters that are most likely to have produced the observations. MLE is discussed in Section 3.6. For ARMA models, this is closely related to least squares.

## Least Squares Estimation

We now discuss least squares for $\operatorname{ARMA}(p, q)$ models via Gauss-Newton. Write the model parameters as $\boldsymbol{\beta}=\left(\phi_{1}, \ldots, \phi_{p}, \theta_{1}, \ldots, \theta_{q}\right)^{\prime}$, and for the ease of discussion, we will put $\mu=0$. Now, write the model in terms of the errors

$$
\begin{equation*}
w_{t}(\boldsymbol{\beta})=x_{t}-\sum_{j=1}^{p} \phi_{j} x_{t-j}-\sum_{k=1}^{q} \theta_{k} w_{t-k}(\boldsymbol{\beta}) \tag{3.31}
\end{equation*}
$$

emphasizing the dependence of the errors on the parameters (recall that $w_{t}=$ $\sum_{j=0}^{\infty} \pi_{j} x_{t-j}$ by invertibilty, and the $\pi_{j}$ are complicated functions of $\boldsymbol{\beta}$ ).

For conditional least squares, we approximate the residual sum of squares by conditioning on $x_{1}, \ldots, x_{p}$ (if $p>0$ ) and $w_{p}=w_{p-1}=w_{p-2}=\cdots=w_{1-q}=0$ (if $q>0$ ), in which case, given $\boldsymbol{\beta}$, we may evaluate (3.31) for $t=p+1, p+2, \ldots, n$. Using this conditioning argument, the conditional error sum of squares is

$$
\begin{equation*}
S_{c}(\boldsymbol{\beta})=\sum_{t=p+1}^{n} w_{t}^{2}(\boldsymbol{\beta}) \tag{3.32}
\end{equation*}
$$

Minimizing $S_{c}(\boldsymbol{\beta})$ with respect to $\boldsymbol{\beta}$ yields the conditional least squares estimates.
If $q=0$, the problem is linear regression and no iterative technique is needed to minimize $S_{c}\left(\phi_{1}, \ldots, \phi_{p}\right)$. For example, for an $\operatorname{AR}(1), x_{t}=\phi x_{t-1}+w_{t}$, the conditional sum of squares is

$$
S_{c}(\phi)=\sum_{t=2}^{n} w_{t}^{2}(\phi)=\sum_{t=2}^{n}\left(x_{t}-\phi x_{t-1}\right)^{2}
$$

Note that we have to start at $t=2$ because $x_{0}$ is not observed. The conditional least squares estimate of $\phi$ follows from simple linear regression wherein,

$$
\widehat{\phi}=\frac{\sum_{t=2}^{n} x_{t} x_{t-1}}{\sum_{t=2}^{n} x_{t-1}^{2}},
$$

which is nearly $\widehat{\rho}(1)$.
If $q>0$, the problem becomes nonlinear regression and we will have to rely on numerical optimization. Gauss-Newton uses an iterative method for solving the problem of minimizing (3.32). We demonstrate the method for an MA(1).

## Example 3.19 Gauss-Newton for an MA(1)

Consider an MA(1) process, $x_{t}=w_{t}+\theta w_{t-1}$. Write the truncated errors as

$$
\begin{equation*}
w_{t}(\theta)=x_{t}-\theta w_{t-1}(\theta), \quad t=1, \ldots, n \tag{3.33}
\end{equation*}
$$

where we condition on $w_{0}(\theta)=0$. Our goal is to find the value of $\theta$ that minimizes $S_{c}(\theta)=\sum_{t=1}^{n} w_{t}^{2}(\theta)$, which is a nonlinear function of $\theta$.

Let $\theta_{(0)}$ be an initial estimate of $\theta$. For example, we could use method of moments. The first-order Taylor expansion ${ }^{3}$ of $w_{t}(\theta)$ at $\theta_{(0)}$ is

$$
\begin{equation*}
w_{t}(\theta) \approx w_{t}\left(\theta_{(0)}\right)-\left(\theta-\theta_{(0)}\right) z_{t}\left(\theta_{(0)}\right) \tag{3.34}
\end{equation*}
$$

where

$$
z_{t}\left(\theta_{(0)}\right)=-\left.\frac{\partial w_{t}(\theta)}{\partial \theta}\right|_{\theta=\theta_{(0)}}
$$

Taking derivatives in (3.33),

$$
\begin{equation*}
\frac{\partial w_{t}(\theta)}{\partial \theta}=-w_{t-1}(\theta)-\theta \frac{\partial w_{t-1}(\theta)}{\partial \theta}, \quad t=1, \ldots, n \tag{3.35}
\end{equation*}
$$

where $\partial w_{0}(\theta) / \partial \theta=0$. Using the notation of (3.34), we can also write (3.35) as

$$
\begin{equation*}
z_{t}(\theta)=w_{t-1}(\theta)-\theta z_{t-1}(\theta), \quad t=1, \ldots, n \tag{3.36}
\end{equation*}
$$

[^10]

Fig. 3.7. ACF and PACF of transformed glacial varves.
where $z_{0}(\theta)=0$. This implies that the derivative sequence is an AR process, which we may easily compute given a value of $\theta$.

The linear approximation of $S_{c}(\theta)$ is found by replacing $w_{t}(\theta)$ by its linear approximation in (3.34),

$$
\begin{equation*}
Q(\theta)=\sum_{t=1}^{n}[\underbrace{w_{t}\left(\theta_{(0)}\right)}_{y_{t}}-\underbrace{\left(\theta-\theta_{(0)}\right)}_{\beta} \underbrace{z_{t}\left(\theta_{(0)}\right)}_{z_{t}}]^{2} \tag{3.37}
\end{equation*}
$$

and this is the quantity that we will minimize. The problem is now simple linear regression (" $y_{t}=\beta z_{t}+\epsilon_{t}$ "), so that

$$
\left(\widehat{\theta-\theta_{(0)}}\right)=\sum_{t=1}^{n} z_{t}\left(\theta_{(0)}\right) w_{t}\left(\theta_{(0)}\right) / \sum_{t=1}^{n} z_{t}^{2}\left(\theta_{(0)}\right)
$$

or

$$
\widehat{\theta}=\theta_{(0)}+\sum_{t=1}^{n} z_{t}\left(\theta_{(0)}\right) w_{t}\left(\theta_{(0)}\right) / \sum_{t=1}^{n} z_{t}^{2}\left(\theta_{(0)}\right)
$$

Consequently, the Gauss-Newton procedure in this case is, on iteration $j+1$, set

$$
\begin{equation*}
\theta_{(j+1)}=\theta_{(j)}+\frac{\sum_{t=1}^{n} z_{t}\left(\theta_{(j)}\right) w_{t}\left(\theta_{(j)}\right)}{\sum_{t=1}^{n} z_{t}^{2}\left(\theta_{(j)}\right)}, \quad j=0,1,2, \ldots, \tag{3.38}
\end{equation*}
$$

where the values in (3.38) are calculated recursively using (3.33) and (3.36). The calculations are stopped when $\left|\theta_{(j+1)}-\theta_{(j)}\right|$, or $\left|Q\left(\theta_{(j+1)}\right)-Q\left(\theta_{(j)}\right)\right|$, are smaller than some preset amount.


Fig. 3.8. Conditional sum of squares versus values of the moving average parameter for the glacial varve example, Example 3.20. Vertical lines indicate the values of the parameter obtained via Gauss-Newton; see Table 3.2 for the actual values.

## Example 3.20 Fitting the Glacial Varve Series

Consider the glacial varve series ( $n=634$ years) analyzed in Example 2.6 and in Problem 2.6. It was argued that a first-order moving average model might fit the logarithmically transformed and differenced series, $\nabla \log \left(x_{t}\right)$, which can be interpreted as being approximately the percentage annual change in the thickness.

The sample ACF and PACF, shown in Figure 3.7, confirm the tendency of $\nabla \log \left(x_{t}\right)$ to behave as a first-order moving average process as the ACF has only a significant peak at lag one and the PACF decreases exponentially. Using Table 3.1, this sample behavior fits that of the MA(1) very well.

Because $\widehat{\rho}(1)=-.4$, it should be clear that $\widehat{\theta}$ is negative. Recalling Example 3.3, the method of moments estimate is found by setting $-.4=\widehat{\theta} /\left(1+\widehat{\theta}^{2}\right)$, which yields two solutions $\widehat{\theta}=-.5$ and -2 . Since -.5 is the invertible solution, we could initialize the Gauss-Newton procedure at this value. However, for most ARMA models, it is not necessary to be so precise as long as the procedure is started in a causal and invertible region. For demonstration purposes, we will start the procedure at $\theta_{(0)}=-.1$.

The results of eleven iterations of the Gauss-Newton procedure, (3.38), are given in Table 3.2. The final estimate is $\widehat{\theta}=\theta_{(11)}=-.773$; interim values and the corresponding value of the conditional sum of squares, $S_{c}(\theta)$ given in (3.32), are also displayed in the table. The final estimate of the error variance is $\widehat{\sigma}_{w}^{2}=$ $148.98 / 632=.236$ with 632 degrees of freedom (one is lost in differencing). The value of the sum of the squared derivatives at convergence is $\sum_{t=1}^{n} z_{t}^{2}\left(\theta_{(11)}\right)=$ 369.73, and consequently, the estimated standard error of $\widehat{\theta}$ is $\sqrt{.236 / 369.73}=$ $.025 ;{ }^{4}$ this leads to a $t$-value of $-.773 / .025=-30.92$ with 632 degrees of freedom.

[^11]Table 3.2. Gauss-Newton Results for Example 3.20

| $j$ | $\theta_{(j)}$ | $S_{c}\left(\theta_{(j)}\right)$ | $\sum_{t=1}^{n} z_{t}^{2}\left(\theta_{(j)}\right)$ |
| :---: | ---: | ---: | :---: |
| 0 | -0.100 | 195.0010 | 183.3464 |
| 1 | -0.250 | 177.7614 | 163.3038 |
| 2 | -0.400 | 165.0027 | 161.6279 |
| 3 | -0.550 | 155.6723 | 182.6432 |
| 4 | -0.684 | 150.2896 | 247.4942 |
| 5 | -0.736 | 149.2283 | 304.3125 |
| 6 | -0.757 | 149.0272 | 337.9200 |
| 7 | -0.766 | 148.9885 | 355.0465 |
| 8 | -0.770 | 148.9812 | 363.2813 |
| 9 | -0.771 | 148.9804 | 365.4045 |
| 10 | -0.772 | 148.9799 | 367.5544 |
| 11 | -0.773 | 148.9799 | 369.7314 |

Figure 3.8 displays the conditional sum of squares, $S_{c}(\theta)$ as a function of $\theta$, as well as indicating the values of each step of the Gauss-Newton algorithm. Note that the Gauss-Newton procedure takes large steps toward the minimum initially, and then takes very small steps as it gets close to the minimizing value. When there is only one parameter, as in this case, it would be easy to evaluate $S_{c}(\theta)$ on a grid of points, and then choose the appropriate value of $\theta$ from the grid search. It would be difficult, however, to perform grid searches when there are many parameters.

In the general case of causal and invertible $\operatorname{ARMA}(p, q)$ models, maximum likelihood estimation and nonlinear least squares estimation (and Yule-Walker estimation in the case of AR models) all lead to optimal estimators.

## Example 3.21 Some Specific Asymptotic Distributions ${ }^{5}$ <br> AR(1):

$$
\begin{equation*}
\widehat{\phi} \sim \mathrm{AN}\left[\phi, n^{-1}\left(1-\phi^{2}\right)\right] . \tag{3.39}
\end{equation*}
$$

AR(2):

$$
\binom{\widehat{\phi}_{1}}{\widehat{\phi}_{2}} \sim \operatorname{AN}\left[\binom{\phi_{1}}{\phi_{2}}, n^{-1}\left(\begin{array}{cc}
1-\phi_{2}^{2} & -\phi_{1}\left(1+\phi_{2}\right)  \tag{3.40}\\
\operatorname{sym} & 1-\phi_{2}^{2}
\end{array}\right)\right] .
$$

MA(1):

$$
\begin{equation*}
\widehat{\theta} \sim \operatorname{AN}\left[\theta, n^{-1}\left(1-\theta^{2}\right)\right] . \tag{3.41}
\end{equation*}
$$

MA(2):

$$
\binom{\widehat{\theta_{1}}}{\widehat{\theta}_{2}} \sim \mathrm{AN}\left[\binom{\theta_{1}}{\theta_{2}}, n^{-1}\left(\begin{array}{cc}
1-\theta_{2}^{2} & \theta_{1}\left(1+\theta_{2}\right)  \tag{3.42}\\
\operatorname{sym} & 1-\theta_{2}^{2}
\end{array}\right)\right] .
$$

ARMA(1,1):

$$
\binom{\widehat{\phi}}{\hat{\theta}} \sim \operatorname{AN}\left[\binom{\phi}{\theta}, n^{-1}\left[\begin{array}{cc}
\left(1-\phi^{2}\right)^{-1} & (1+\phi \theta)^{-1}  \tag{3.43}\\
\operatorname{sym} & \left(1-\theta^{2}\right)^{-1}
\end{array}\right]^{-1}\right] .
$$

[^12]
## Example 3.22 Overfitting Caveat

The asymptotic behavior of the parameter estimators gives us an additional insight into the problem of fitting ARMA models to data. For example, suppose a time series follows an $\operatorname{AR}(1)$ process and we decide to fit an $\operatorname{AR}(2)$ to the data. Do any problems occur in doing this? More generally, why not simply fit large-order AR models to make sure that we capture the dynamics of the process? After all, if the process is truly an $\operatorname{AR}(1)$, the other autoregressive parameters will not be significant. The answer is that if we overfit, we obtain less efficient, or less precise parameter estimates. For example, if we fit an $\operatorname{AR}(1)$ to an $\operatorname{AR}(1)$ process, for large $n, \operatorname{var}\left(\widehat{\phi}_{1}\right) \approx n^{-1}\left(1-\phi_{1}^{2}\right)$. But, if we fit an $\operatorname{AR}(2)$ to the $\operatorname{AR}(1)$ process, for large $n$, $\operatorname{var}\left(\widehat{\phi}_{1}\right) \approx n^{-1}\left(1-\phi_{2}^{2}\right)=n^{-1}$ because $\phi_{2}=0$. Thus, the variance of $\phi_{1}$ has been inflated, making the estimator less precise.

We do want to mention, however, that overfitting can be used as a diagnostic tool. For example, if we fit an $\operatorname{AR}(2)$ model to the data and are satisfied with that model, then adding one more parameter and fitting an $\operatorname{AR}(3)$ should lead to approximately the same model as in the $\operatorname{AR}(2)$ fit. We will discuss model diagnostics in more detail in §3.8.

### 3.5 Forecasting

In forecasting, the goal is to predict future values of a time series, $x_{n+m}, m=1,2, \ldots$, based on the data collected to the present. Throughout this section, we will assume $x_{t}$ is stationary and the model parameters are known. When the parameters are unknown, we replace them with their estimates.

First, consider one-step-ahead prediction. That is, given $\left\{x_{1}, \ldots, x_{n}\right\}$, we wish to forecast the value of the time series at the next time point, $x_{n+1}$. The best linear predictor (BLP) of $x_{n+1}$ is of the form

$$
\begin{equation*}
x_{n+1}^{n}=\phi_{n 1} x_{n}+\phi_{n 2} x_{n-1}+\cdots+\phi_{n n} x_{1} \tag{3.44}
\end{equation*}
$$

and we must solve for the coefficients to minimize the mean squared prediction error (MSPE) given by

$$
P_{n+1}^{n}=E\left[x_{n+1}-x_{n+1}^{n}\right]^{2}=E\left[x_{n+1}-\left(\phi_{n 1} x_{n}+\cdots+\phi_{n n} x_{1}\right)\right]^{2}
$$

Generally, the coefficients depend on the sample size $n$, and this is made explicit. The second subscript indicates the distance between the future and the data values. For example, $\phi_{n j}$ is the coefficient of $x_{n+1-j}$ because it is $j$ time periods away from $x_{n+1}$.

The coefficients $\left\{\phi_{n 1}, \phi_{n 2}, \ldots, \phi_{n n}\right\}$ satisfy

$$
\begin{equation*}
\sum_{j=1}^{n} \phi_{n j} \gamma(k-j)=\gamma(k), \quad k=1, \ldots, n \tag{3.45}
\end{equation*}
$$

which is similar to the normal equations obtained in least squares regression. If $n$ is large (which is typical in time series), solving (3.45) directly is computationally
prohibitive. There are, however, iterative solutions that do not require any matrix inversion. In particular, we mention the recursive solution due to Levinson (1947) and Durbin (1960).

## Property 3.3 The Durbin-Levinson Algorithm

The coefficients in (3.44) and the corresponding MSPE can be solved iteratively as follows:

$$
\begin{equation*}
\phi_{00}=0, \quad P_{1}^{0}=\gamma(0) \tag{3.46}
\end{equation*}
$$

For $n \geq 1$,

$$
\begin{equation*}
\phi_{n n}=\frac{\rho(n)-\sum_{k=1}^{n-1} \phi_{n-1, k} \rho(n-k)}{1-\sum_{k=1}^{n-1} \phi_{n-1, k} \rho(k)}, \quad P_{n+1}^{n}=P_{n}^{n-1}\left(1-\phi_{n n}^{2}\right), \tag{3.47}
\end{equation*}
$$

where, for $n \geq 2, \phi_{n k}=\phi_{n-1, k}-\phi_{n n} \phi_{n-1, n-k}, k=1,2, \ldots, n-1$.
An important consequence of the Durbin-Levinson algorithm is as follows.

## Property 3.4 Iterative Solution for the PACF

The PACF of a stationary process $x_{t}$, can be obtained iteratively via (3.47) as $\phi_{n n}$, for $n=1,2, \ldots$.

Property 3.4 follows from the fact that, in (3.44), $\phi_{n n}$ is the (auto)regression coefficient of $x_{n+1}$ on $x_{1}$, and so it measures the effect of $x_{1}$ on $x_{n+1}$ with everything in the middle being held constant. The fact that we are regressing $x_{t}$ on itself, makes the regression coefficient $\phi_{n n}$ a partial correlation.

We note that given data from a stationary process, we could substitute $\widehat{\gamma}(k)$ for $\gamma(k)$ in (3.45) and then solve for the $\phi$ s to obtain their estimates and consequently the sample-based forecasts and the sample PACF values (without having to run repeated regressions). For AR models, forecasting simplifies.

## Example 3.23 Prediction for an AR(2)

Suppose we have a causal $\operatorname{AR}(2)$ process $x_{t}=\phi_{1} x_{t-1}+\phi_{2} x_{t-2}+w_{t}$, and one observation $x_{1}$. Then, using equation (3.45), $\phi_{11} \gamma(0)=\gamma(1)$, the one-step-ahead prediction of $x_{2}$ based on $x_{1}$ is

$$
x_{2}^{1}=\phi_{11} x_{1}=\frac{\gamma(1)}{\gamma(0)} x_{1}=\rho(1) x_{1}
$$

Now, suppose we want the one-step-ahead prediction of $x_{3}$ based on two observations $x_{1}$ and $x_{2}$; i.e., $x_{3}^{2}=\phi_{21} x_{2}+\phi_{22} x_{1}$. We could use (3.45)

$$
\begin{aligned}
& \phi_{21} \gamma(0)+\phi_{22} \gamma(1)=\gamma(1) \\
& \phi_{21} \gamma(1)+\phi_{22} \gamma(0)=\gamma(2)
\end{aligned}
$$

to solve for $\phi_{21}$ and $\phi_{22}$, but it should be apparent from the model that $x_{3}^{2}=$ $\phi_{1} x_{2}+\phi_{2} x_{1}$. It is easy to verify that, for $n \geq 2$,

$$
x_{n+1}^{n}=\phi_{1} x_{n}+\phi_{2} x_{n-1} .
$$

That is, $\phi_{n 1}=\phi_{1}, \phi_{n 2}=\phi_{2}$, and $\phi_{n j}=0$, for $j=3,4, \ldots, n$.

In general, if the time series is an $\operatorname{AR}(p)$ process, then, for $n \geq p$,

$$
\begin{equation*}
x_{n+1}^{n}=\phi_{1} x_{n}+\phi_{2} x_{n-1}+\cdots+\phi_{p} x_{n-p+1} . \tag{3.48}
\end{equation*}
$$

So far, we have concentrated on one-step-ahead prediction, but $m$-step-ahead prediction is similar. The BLP of $x_{n+m}$ for any $m \geq 1$ given data $\left\{x_{1}, \ldots, x_{n}\right\}$ yields the $m$-step-ahead predictor

$$
\begin{equation*}
x_{n+m}^{n}=\phi_{n 1}^{(m)} x_{n}+\phi_{n 2}^{(m)} x_{n-1}+\cdots+\phi_{n n}^{(m)} x_{1}, \tag{3.49}
\end{equation*}
$$

where $\left\{\phi_{n 1}^{(m)}, \phi_{n 2}^{(m)}, \ldots, \phi_{n n}^{(m)}\right\}$ satisfy

$$
\begin{equation*}
\sum_{j=1}^{n} \phi_{n j}^{(m)} \gamma(k-j)=\gamma(m+k-1), \quad k=1, \ldots, n \tag{3.50}
\end{equation*}
$$

The mean square m-step-ahead prediction error is $P_{n+m}^{n}=E\left(x_{n+m}-x_{n+m}^{n}\right)^{2}$.

## Forecasting ARMA Processes

Because ARMA models are invertible; i.e., $w_{t}=x_{t}+\sum_{j=1}^{\infty} \pi_{j} x_{t-j}$, we may write

$$
x_{n+m}=-\sum_{j=1}^{\infty} \pi_{j} x_{n+m-j}+w_{n+m} .
$$

If we had the infinite history $\left\{x_{n}, x_{n-1}, \ldots, x_{1}, x_{0}, x_{-1}, \ldots\right\}$, of the data available, we would predict $x_{n+m}$ by

$$
x_{n+m}^{n}=-\sum_{j=1}^{\infty} \pi_{j} x_{n+m-j}^{n},
$$

successively for $m=1,2, \ldots$ In this case, $x_{t}^{n}=x_{t}$ for $t=n, n-1, \ldots$ We only have the actual data $\left\{x_{n}, x_{n-1}, \ldots, x_{1}\right\}$ available, but a practical solution is to truncate the forecasts as

$$
x_{n+m}^{n}=-\sum_{j=1}^{n+m-1} \pi_{j} x_{n+m-j}^{n},
$$

with $x_{t}^{n}=x_{t}$ for $1 \leq t \leq n$. As displayed in (3.48), for $\operatorname{AR}(p)$ models, truncated forecasts are exact when $n>p$. For ARMA models in general, as long as $n$ is large, the approximation works well because the $\pi$-weights are going to zero exponentially fast. For ARMA models, truncated forecasts can be obtained fairly simply as discussed in the following example.

## Example 3.24 Forecasting an ARMA $(1,1)$ Series

Given data $x_{1}, \ldots, x_{n}$, for forecasting purposes, write the model as

$$
x_{n+1}=\phi x_{n}+w_{n+1}+\theta w_{n} .
$$

Note that $w_{n}$ is the current error, which is a function of the data. But $w_{n+1}$ is a future error that is not a function of the data. The one-step-ahead truncated forecast is

$$
x_{n+1}^{n}=\phi x_{n}+0+\theta w_{n}^{n},
$$

where $w_{n}^{n}$ is the estimate of $w_{n}$ given the data. For $m \geq 2$, we have

$$
x_{n+m}^{n}=\phi x_{n+m-1}^{n},
$$

which can be calculated recursively, $m=2,3, \ldots$.
To calculate $w_{n}^{n}$, which is needed to initialize the successive forecasts, the model can be written as $w_{t}=x_{t}-\phi x_{t-1}-\theta w_{t-1}$ for $t=1, \ldots, n$. For truncated forecasting put $w_{0}^{n}=0, x_{0}=0$, and then iterate the errors forward in time

$$
w_{t}^{n}=x_{t}-\phi x_{t-1}-\theta w_{t-1}^{n}, \quad t=1, \ldots, n .
$$

For large $n$, it can be shown that the mean squared prediction error for $\operatorname{ARMA}(p, q)$ models is approximately ${ }^{6}$

$$
\begin{equation*}
P_{n+m}^{n}=\sigma_{w}^{2} \sum_{j=0}^{m-1} \psi_{j}^{2} \tag{3.51}
\end{equation*}
$$

To assess the precision of the forecasts, prediction intervals are typically calculated along with the forecasts. In general, $(1-\alpha)$ prediction intervals are of the form

$$
\begin{equation*}
x_{n+m}^{n} \pm c_{\frac{\alpha}{2}} \sqrt{P_{n+m}^{n}} \tag{3.52}
\end{equation*}
$$

where $c_{\alpha / 2}$ is chosen to get the desired degree of confidence. For example, if the process is Gaussian, then choosing $c_{\alpha / 2}=2$ will yield an approximate $95 \%$ prediction interval for $x_{n+m}$. If we are interested in establishing prediction intervals over more than one time period, then $c_{\alpha / 2}$ should be adjusted appropriately, for example, by using Bonferroni's inequality [see Johnson and Wichern, 1992, Chapter 5].

## Example 3.25 Forecasting the Recruitment Series

In Example 3.15 we fit an $\operatorname{AR}(2)$ model to the Recruitment series using OLS. Here, we use MLE:


[^13]

Fig. 3.9. Twenty-four month forecasts for the Recruitment series. The actual data shown are from about January 1980 to September 1987, and then the forecasts plus and minus one standard error are displayed.

The results are nearly the same as using OLS. Using the parameter estimates as the actual parameter values, the forecasts are calculated as

$$
x_{n+m}^{n}=6.80+1.35 x_{n+m-1}^{n}-.46 x_{n+m-2}^{n}
$$

for $n=453$ and $m=1,2, \ldots, 12$. Recall that $x_{t}^{s}=x_{t}$ when $t \leq s$. The forecasts errors $P_{n+m}^{n}$ are calculated using (3.51). Here, $\widehat{\sigma}_{w}^{2}=89.33$, and the $\psi$-weights satisfy $\psi_{j}=1.35 \psi_{j-1}-.46 \psi_{j-2}$ for $j \geq 2$, where $\psi_{0}=1$ and $\psi_{1}=1.35$. Thus, for $n=453$,

$$
P_{n+1}^{n}=89.33, P_{n+2}^{n}=89.33\left(1+1.35^{2}\right), P_{n+3}^{n}=89.33\left(1+1.35^{2}+\left[1.35^{2}-.46\right]^{2}\right),
$$

and so on.
Figure 3.9 shows the result of forecasting the Recruitment series over a 24-month horizon, $m=1,2, \ldots, 24$, obtained in R as sarima. for (rec, 24, 2, 0, 0)
Note how the forecast levels off to the mean quickly and the prediction intervals are wide and become constant. That is, because of the short memory, the forecasts settle to the mean, $\mu_{x}$, of the process, and the MSPE becomes $\gamma_{x}(0)=\operatorname{var}\left(x_{t}\right)$.

We complete this section with a brief discussion of backcasting. In backcasting, we want to predict $x_{1-m}$, for $m=1,2, \ldots$, based on the data $\left\{x_{1}, \ldots, x_{n}\right\}$. Write the backcast as

$$
\begin{equation*}
x_{1-m}^{n}=\sum_{j=1}^{n} \alpha_{n j}^{(m)} x_{j} \tag{3.53}
\end{equation*}
$$

Analogous to (3.50), the prediction equations (assuming $\mu_{x}=0$ ) are

$$
\begin{equation*}
\sum_{j=1}^{n} \alpha_{n j}^{(m)} \gamma(k-j)=\gamma(m+k-1), \quad k=1, \ldots, n \tag{3.54}
\end{equation*}
$$



Fig. 3.10. Display for Example 3.26; backcasts from a simulated ARMA(1, 1).

These equations are precisely the prediction equations for forward prediction. That is, $\alpha_{n j}^{(m)} \equiv \phi_{n j}^{(m)}$, for $j=1, \ldots, n$, where the $\phi_{n j}^{(m)}$ are given by (3.50). Finally, the backcasts are given by

$$
\begin{equation*}
x_{1-m}^{n}=\phi_{n 1}^{(m)} x_{1}+\cdots+\phi_{n n}^{(m)} x_{n}, \quad m=1,2, \ldots \tag{3.55}
\end{equation*}
$$

The implication of this result is that to backcast, you simply reverse the order of the data and forecast the reversed data set.

## Example 3.26 Backcasting an ARMA Process

In the following R code, we backcast an $\operatorname{ARMA}(1,1)$ by reversing simulated data, fitting the model and then predicting the reversed data; see Figure 3.10.
set.seed (90210)
$\mathrm{x}=\operatorname{arima} . \operatorname{sim}(\operatorname{list}($ order $=c(1,0,1)$, $\mathrm{ar}=.9, \operatorname{ma}=.5), \mathrm{n}=100)$
$\mathrm{xr}=$ replace(x, TRUE, $\mathrm{rev}(\mathrm{x})$ ) \# xr is the reversed data
$\mathrm{pxr}=$ predict $(\operatorname{arima}(\mathrm{xr}, \operatorname{order}=\mathrm{c}(1,0,1)), 10)$ \# predict the reversed data
pxrp $=\operatorname{rev}(p x r \$ p r e d) \quad$ \# reorder the predictors (for plotting)
pxrse $=\operatorname{rev}($ pxr\$se $)$ \# reorder the SEs
$\mathrm{nx}=\mathrm{ts}(\mathrm{c}(\mathrm{pxrp}, \mathrm{x})$, start=-9) \# attach the backcasts to the data
plot(nx, ylab=expression(X[~t]), main='Backcasting')
lines (-9:0, nx[1:10], col=2, type='o', pch=20)
lines $(-9: 0, \mathrm{nx}[1: 10]+\mathrm{pxrse}, \mathrm{col}=4, \mathrm{lty}=2$ )
lines (-9:0, nx[1:10] - pxrse, col=4, lty=2)

### 3.6 Maximum Likelihood Estimation **

For a normal $\operatorname{ARMA}(p, q)$ model, the optimal way to estimate the parameters is to use either maximum likelihood estimation, or unconditional least squares estimation. Without going into general details, we describe the technique for an $\operatorname{AR}(1)$ model.

[^14]
## Example 3.27 Estimation for an AR(1)

Let

$$
\begin{equation*}
x_{t}=\mu+\phi\left(x_{t-1}-\mu\right)+w_{t} \tag{3.56}
\end{equation*}
$$

where $|\phi|<1$ and $w_{t} \sim \operatorname{iid} \mathrm{~N}\left(0, \sigma_{w}^{2}\right)$. Given data $x_{1}, x_{2}, \ldots, x_{n}$, we seek the likelihood

$$
L\left(\mu, \phi, \sigma_{w}^{2}\right)=f\left(x_{1}, x_{2}, \ldots, x_{n} \mid \mu, \phi, \sigma_{w}^{2}\right)
$$

In the case of an $\operatorname{AR}(1)$, we may write the likelihood as

$$
L\left(\mu, \phi, \sigma_{w}^{2}\right)=f\left(x_{1}\right) f\left(x_{2} \mid x_{1}\right) \cdots f\left(x_{n} \mid x_{n-1}\right)
$$

where we have dropped the parameters in the densities, $f(\cdot)$, to ease the notation. Because, for $t>1, x_{t} \mid x_{t-1} \sim \mathrm{~N}\left(\mu+\phi\left(x_{t-1}-\mu\right), \sigma_{w}^{2}\right)$, we have

$$
f\left(x_{t} \mid x_{t-1}\right)=f_{w}\left[\left(x_{t}-\mu\right)-\phi\left(x_{t-1}-\mu\right)\right]
$$

where $f_{w}(\cdot)$ is the density of $w_{t}$, that is, the normal density with mean zero and variance $\sigma_{w}^{2}$. We may then write the likelihood as

$$
L\left(\mu, \phi, \sigma_{w}\right)=f\left(x_{1}\right) \prod_{t=2}^{n} f_{w}\left[\left(x_{t}-\mu\right)-\phi\left(x_{t-1}-\mu\right)\right] .
$$

To find $f\left(x_{1}\right)$, we can use the causal representation

$$
x_{1}=\mu+\sum_{j=0}^{\infty} \phi^{j} w_{1-j}
$$

to see that $x_{1}$ is normal, with mean $\mu$ and variance $\sigma_{w}^{2} /\left(1-\phi^{2}\right)$. Finally, for an $\operatorname{AR}(1)$, the likelihood is

$$
\begin{equation*}
L\left(\mu, \phi, \sigma_{w}^{2}\right)=\left(2 \pi \sigma_{w}^{2}\right)^{-n / 2}\left(1-\phi^{2}\right)^{1 / 2} \exp \left[-\frac{S(\mu, \phi)}{2 \sigma_{w}^{2}}\right] \tag{3.57}
\end{equation*}
$$

where

$$
\begin{equation*}
S(\mu, \phi)=\left(1-\phi^{2}\right)\left(x_{1}-\mu\right)^{2}+\sum_{t=2}^{n}\left[\left(x_{t}-\mu\right)-\phi\left(x_{t-1}-\mu\right)\right]^{2} . \tag{3.58}
\end{equation*}
$$

Typically, $S(\mu, \phi)$ is called the unconditional sum of squares. We could have also considered the estimation of $\mu$ and $\phi$ using unconditional least squares, that is, estimation by minimizing $S(\mu, \phi)$.

Taking the partial derivative of the $\log$ of (3.57) with respect to $\sigma_{w}^{2}$ and setting the result equal to zero, we see that for any given values of $\mu$ and $\phi$ in the parameter space, $\sigma_{w}^{2}=n^{-1} S(\mu, \phi)$ maximizes the likelihood. Thus, the maximum likelihood estimate of $\sigma_{w}^{2}$ is

$$
\begin{equation*}
\widehat{\sigma}_{w}^{2}=n^{-1} S(\widehat{\mu}, \widehat{\phi}), \tag{3.59}
\end{equation*}
$$

where $\widehat{\mu}$ and $\widehat{\phi}$ are the MLEs of $\mu$ and $\phi$, respectively. If we replace $n$ in (3.59) by $n-2$, we would obtain the unconditional least squares estimate of $\sigma_{w}^{2}$.

If, in (3.57), we take logs, replace $\sigma_{w}^{2}$ by $\widehat{\sigma}_{w}^{2}$, and ignore constants, $\widehat{\mu}$ and $\widehat{\phi}$ are the values that minimize the criterion function

$$
\begin{equation*}
l(\mu, \phi)=\log \left[n^{-1} S(\mu, \phi)\right]-n^{-1} \log \left(1-\phi^{2}\right) \tag{3.60}
\end{equation*}
$$

that is, $l(\mu, \phi) \propto-2 \log L\left(\mu, \phi, \widehat{\sigma}_{w}^{2}\right) .^{7}$ Because (3.58) and (3.60) are complicated functions of the parameters, the minimization of $l(\mu, \phi)$ or $S(\mu, \phi)$ is accomplished numerically. In the case of AR models, we have the advantage that, conditional on initial values, they are linear models. That is, we can drop the term in the likelihood that causes the nonlinearity. Conditioning on $x_{1}$, the conditional likelihood becomes

$$
\begin{align*}
L\left(\mu, \phi, \sigma_{w}^{2} \mid x_{1}\right) & =\prod_{t=2}^{n} f_{w}\left[\left(x_{t}-\mu\right)-\phi\left(x_{t-1}-\mu\right)\right] \\
& =\left(2 \pi \sigma_{w}^{2}\right)^{-(n-1) / 2} \exp \left[-\frac{S_{c}(\mu, \phi)}{2 \sigma_{w}^{2}}\right] \tag{3.61}
\end{align*}
$$

where the conditional sum of squares is

$$
\begin{equation*}
S_{c}(\mu, \phi)=\sum_{t=2}^{n}\left[\left(x_{t}-\mu\right)-\phi\left(x_{t-1}-\mu\right)\right]^{2} \tag{3.62}
\end{equation*}
$$

The conditional MLE of $\sigma_{w}^{2}$ is

$$
\begin{equation*}
\widehat{\sigma}_{w}^{2}=S_{c}(\widehat{\mu}, \widehat{\phi}) /(n-1) \tag{3.63}
\end{equation*}
$$

and $\widehat{\mu}$ and $\widehat{\phi}$ are the values that minimize the conditional sum of squares, $S_{c}(\mu, \phi)$. Letting $\alpha=\mu(1-\phi)$, the conditional sum of squares can be written as

$$
\begin{equation*}
S_{c}(\mu, \phi)=\sum_{t=2}^{n}\left[x_{t}-\left(\alpha+\phi x_{t-1}\right)\right]^{2} \tag{3.64}
\end{equation*}
$$

The problem is now the linear regression problem stated in §2.2. Following the results from least squares estimation, we have $\widehat{\alpha}=\bar{x}_{(2)}-\widehat{\phi} \bar{x}_{(1)}$, where $\bar{x}_{(1)}=$ $(n-1)^{-1} \sum_{t=1}^{n-1} x_{t}$, and $\bar{x}_{(2)}=(n-1)^{-1} \sum_{t=2}^{n} x_{t}$, and the conditional estimates are then

$$
\begin{gather*}
\widehat{\mu}=\frac{\bar{x}_{(2)}-\widehat{\phi} \bar{x}_{(1)}}{1-\widehat{\phi}}  \tag{3.65}\\
\widehat{\phi}=\frac{\sum_{t=2}^{n}\left(x_{t}-\bar{x}_{(2)}\right)\left(x_{t-1}-\bar{x}_{(1)}\right)}{\sum_{t=2}^{n}\left(x_{t-1}-\bar{x}_{(1)}\right)^{2}} . \tag{3.66}
\end{gather*}
$$

From (3.65) and (3.66), we see that $\widehat{\mu} \approx \bar{x}$ and $\widehat{\phi} \approx \widehat{\rho}(1)$. That is, the Yule-Walker estimators and the conditional least squares estimators are approximately the same.

[^15]The only difference is the inclusion or exclusion of terms involving the endpoints, $x_{1}$ and $x_{n}$. We can also adjust the estimate of $\sigma_{w}^{2}$ in (3.63) to be equivalent to the least squares estimator, that is, divide $S_{c}(\widehat{\mu}, \widehat{\phi})$ by $(n-3)$ instead of $(n-1)$ in (3.63).

For general $\operatorname{AR}(p)$ models, maximum likelihood estimation, unconditional least squares, and conditional least squares follow analogously to the AR(1) example. For ARMA models in general, the densities $f\left(x_{t} \mid x_{1}, \ldots, x_{t-1}\right)$ that form the likelihood are obtained using the forecasting methods discussed in the previous section; details are in the blue version of the text.

### 3.7 Integrated Models

In Chapters 1 and 2, we saw that if $x_{t}$ is a random walk, $x_{t}=x_{t-1}+w_{t}$, then by differencing $x_{t}$, we find that $\nabla x_{t}=w_{t}$ is stationary. In many situations, time series can be thought of as being composed of two components, a nonstationary trend component and a zero-mean stationary component. For example, in §2.2 we considered the model

$$
\begin{equation*}
x_{t}=\mu_{t}+y_{t} \tag{3.67}
\end{equation*}
$$

where $\mu_{t}=\beta_{0}+\beta_{1} t$ and $y_{t}$ is stationary. Differencing such a process will lead to a stationary process:

$$
\nabla x_{t}=x_{t}-x_{t-1}=\beta_{1}+y_{t}-y_{t-1}=\beta_{1}+\nabla y_{t} .
$$

Another model that leads to first differencing is the case in which $\mu_{t}$ in (3.67) is stochastic and slowly varying according to a random walk. That is,

$$
\mu_{t}=\mu_{t-1}+v_{t}
$$

where $v_{t}$ is stationary. In this case,

$$
\nabla x_{t}=v_{t}+\nabla y_{t}
$$

is stationary. If $\mu_{t}$ in (3.67) is quadratic, $\mu_{t}=\beta_{0}+\beta_{1} t+\beta_{2} t^{2}$, then the differenced series $\nabla^{2} y_{t}$ is stationary. Stochastic trend models can also lead to higher order differencing. For example, suppose

$$
\mu_{t}=\mu_{t-1}+v_{t} \quad \text { and } \quad v_{t}=v_{t-1}+e_{t},
$$

where $e_{t}$ is stationary. Then, $\nabla x_{t}=v_{t}+\nabla y_{t}$ is not stationary, but

$$
\nabla^{2} x_{t}=e_{t}+\nabla^{2} y_{t}
$$

is stationary.
The integrated ARMA, or ARIMA, model is a broadening of the class of ARMA models to include differencing. The basic idea is that if differencing the data at some order $d$ produces an ARMA process, then the original process is said to be ARIMA.

Definition 3.7 A process $x_{t}$ is said to be $\operatorname{ARIMA}(p, d, q)$ if

$$
\nabla^{d} x_{t}=(1-B)^{d} x_{t}
$$

is $\operatorname{ARMA}(p, q)$. In general, we will write the model as

$$
\begin{equation*}
\phi(B)(1-B)^{d} x_{t}=\theta(B) w_{t} . \tag{3.68}
\end{equation*}
$$

If $E\left(\nabla^{d} x_{t}\right)=\mu$, we write the model as

$$
\phi(B)(1-B)^{d} x_{t}=\delta+\theta(B) w_{t}
$$

where $\delta=\mu\left(1-\phi_{1}-\cdots-\phi_{p}\right)$.
It should be clear that, since $y_{t}=\nabla^{d} x_{t}$ is ARMA, we can use $\S 3.5$ methods to obtain forecasts of $y_{t}$, which in turn lead to forecasts for $x_{t}$. For example, if $d=1$, given forecasts $y_{n+m}^{n}$ for $m=1,2, \ldots$, we have $y_{n+m}^{n}=x_{n+m}^{n}-x_{n+m-1}^{n}$, so that

$$
x_{n+m}^{n}=y_{n+m}^{n}+x_{n+m-1}^{n}
$$

with initial condition $x_{n+1}^{n}=y_{n+1}^{n}+x_{n}\left(\right.$ noting $\left.x_{n}^{n}=x_{n}\right)$.
It is a little more difficult to obtain the prediction errors $P_{n+m}^{n}$, but for large $n$, the approximation used in $\S 3.5$, equation (3.51), works well. That is, the mean-squared prediction error can be approximated by

$$
\begin{equation*}
P_{n+m}^{n}=\sigma_{w}^{2} \sum_{j=0}^{m-1} \psi_{j}^{* 2}, \tag{3.69}
\end{equation*}
$$

where $\psi_{j}^{*}$ is the coefficient of $z^{j}$ in $\psi^{*}(z)=\theta(z) / \phi(z)(1-z)^{d}$.
To better understand forecasting integrated models, we examine the properties of some simple cases.

## Example 3.28 Random Walk with Drift

To fix ideas, we begin by considering the random walk with drift model first presented in Example 1.9, that is,

$$
x_{t}=\delta+x_{t-1}+w_{t},
$$

for $t=1,2, \ldots$, and $x_{0}=0$. Technically, the model is not ARIMA, but we could include it trivially as an $\operatorname{ARIMA}(0,1,0)$ model. Given data $x_{1}, \ldots, x_{n}$, the one-step-ahead forecast is given by

$$
x_{n+1}^{n}=E\left(x_{n+1} \mid x_{n}, \ldots, x_{1}\right)=E\left(\delta+x_{n}+w_{n+1} \mid x_{n}, \ldots, x_{1}\right)=\delta+x_{n} .
$$

The two-step-ahead forecast is given by $x_{n+2}^{n}=\delta+x_{n+1}^{n}=2 \delta+x_{n}$, and consequently, the $m$-step-ahead forecast, for $m=1,2, \ldots$, is

$$
\begin{equation*}
x_{n+m}^{n}=m \delta+x_{n}, \tag{3.70}
\end{equation*}
$$

To obtain the forecast errors, it is convenient to recall equation (1.4), i.e., $x_{n}=n \delta+\sum_{j=1}^{n} w_{j}$, in which case we may write

$$
x_{n+m}=(n+m) \delta+\sum_{j=1}^{n+m} w_{j}=m \delta+x_{n}+\sum_{j=n+1}^{n+m} w_{j} .
$$

From this it follows that the $m$-step-ahead prediction error is given by

$$
\begin{equation*}
P_{n+m}^{n}=E\left(x_{n+m}-x_{n+m}^{n}\right)^{2}=E\left(\sum_{j=n+1}^{n+m} w_{j}\right)^{2}=m \sigma_{w}^{2} \tag{3.71}
\end{equation*}
$$

Unlike the stationary case, as the forecast horizon grows, the prediction errors, (3.71), increase without bound and the forecasts follow a straight line with slope $\delta$ emanating from $x_{n}$.

## Example 3.29 IMA (1, 1) and EWMA

The $\operatorname{ARIMA}(0,1,1)$, or $\operatorname{IMA}(1,1)$ model is of interest because many economic time series can be successfully modeled this way. The model leads to a frequently used forecasting method called exponentially weighted moving averages (EWMA). We will write the model as

$$
\begin{equation*}
x_{t}=x_{t-1}+w_{t}-\lambda w_{t-1} \tag{3.72}
\end{equation*}
$$

with $|\lambda|<1$, for $t=1,2, \ldots$, and $x_{0}=0$, because this model formulation is easier to work with here, and it leads to the standard representation for EWMA. We could have included a drift term in (3.72), as was done in the previous example, but for the sake of simplicity, we leave it out of the discussion. If we write

$$
y_{t}=w_{t}-\lambda w_{t-1},
$$

we may write (3.72) as $x_{t}=x_{t-1}+y_{t}$. Because $|\lambda|<1, y_{t}$ has an invertible representation, $y_{t}+\sum_{j=1}^{\infty} \lambda^{j} y_{t-j}=w_{t}$, and substituting $y_{t}=x_{t}-x_{t-1}$, we may write

$$
\begin{equation*}
x_{t}=\sum_{j=1}^{\infty}(1-\lambda) \lambda^{j-1} x_{t-j}+w_{t} \tag{3.73}
\end{equation*}
$$

as an approximation for large $t$ (put $x_{t}=0$ for $t \leq 0$ ). Verification of (3.73) is left to the reader (Problem 3.11). Using the approximation (3.73), we have that the approximate one-step-ahead predictor is

$$
\begin{equation*}
x_{n+1}^{n}=(1-\lambda) x_{n}+\lambda x_{n}^{n-1}, \tag{3.74}
\end{equation*}
$$

because $x_{n}^{n-1}=\sum_{j=1}^{\infty}(1-\lambda) \lambda^{j-1} x_{n-j}$ and $w_{n+1}^{n}=0$. From (3.74), we see that the new forecast is a linear combination of the old forecast and the new observation. The mean-square prediction error can be approximated using (3.69) by noting that $\psi^{*}(z)=(1-\lambda z) /(1-z)=1+(1-\lambda) \sum_{j=1}^{\infty} z^{j}$ for $|z|<1$; consequently, for large $n$, (3.69) leads to


Fig. 3.11. Output for Example 3.29: Simulated data with an EWMA superimposed.

$$
P_{n+m}^{n} \approx \sigma_{w}^{2}\left[1+(m-1)(1-\lambda)^{2}\right] .
$$

In EWMA, the parameter $1-\lambda$ is often called the smoothing parameter and is restricted to be between zero and one. Larger values of $\lambda$ lead to smoother forecasts. This method of forecasting is popular because it is easy to use; we need only retain the previous forecast value and the current observation to forecast the next time period. In the following, we show how to generate 100 observations from an IMA( 1,1 ) model with $\lambda=-\theta=.8$ and then calculate and display the fitted EWMA superimposed on the data. This is accomplished using the Holt-Winters command in R (see the help file ?HoltWinters for details). The results are displayed in Figure 3.11

```
set.seed(666)
x = arima.sim(list(order = c(0,1,1), ma = -0.8), n = 100)
(x.ima = HoltWinters(x, beta=FALSE, gamma=FALSE)) # \alpha below is 1-\lambda
    Smoothing parameter: alpha: 0.1663072
plot(x.ima, main='EWMA')
```


### 3.8 Building ARIMA Models

There are a few basic steps to fitting ARIMA models to time series data. These steps involve (i) plotting the data, (ii) possibly transforming the data, (iii) identifying the dependence orders of the model, (iv) parameter estimation, (v) diagnostics, (vi) and model choice. First, as with any data analysis, we should construct a time plot of the data, and inspect the graph for any anomalies. If, for example, the variability in the data grows with time, it will be necessary to transform the data to stabilize the variance. In such cases, the Box-Cox class of power transformations, equation (2.35), could be employed. Also, the particular application might suggest an appropriate transformation. For example, suppose a process evolves as a fairly small and stable percent-change, such as an investment. Then a reasonable model is

$$
x_{t}=\left(1+p_{t}\right) x_{t-1},
$$

where $x_{t}$ is the value of the investment at time $t$ and $p_{t}$ is the percentage-change from period $t-1$ to $t$, which may be negative. Recall Figure 1.2 where $p_{t} \equiv .05$ and was compared to Figure 1.1. Taking logs we have


Fig. 3.12. Quarterly U.S. GNP from 1947(1) to 2002(3).

$$
\log \left(x_{t}\right)=\log \left(1+p_{t}\right)+\log \left(x_{t-1}\right)
$$

or

$$
\nabla \log \left(x_{t}\right)=\log \left(1+p_{t}\right)
$$

If the percent change $p_{t}$ stays relatively small in magnitude, then $\log \left(1+p_{t}\right) \approx p_{t}{ }^{8}$ and, thus,

$$
\nabla \log \left(x_{t}\right) \approx p_{t}
$$

will be a relatively stable process. Frequently, $\nabla \log \left(x_{t}\right)$ is called the return or growth rate. This general idea was used in Example 3.20, and we will use it again in Example 3.30.

After suitably transforming the data, the next step is to identify preliminary values of the autoregressive order, $p$, the order of differencing, $d$, and the moving average order, $q$. We have already addressed, in part, the problem of selecting $d$. A time plot of the data will typically suggest whether any differencing is needed. If differencing is called for, then difference the data once, $d=1$, and inspect the time plot of $\nabla x_{t}$. If additional differencing is necessary, then try differencing again and inspect a time plot of $\nabla^{2} x_{t}$. Be careful not to overdifference because this may introduce dependence where none exists. In addition to time plots, the sample ACF can help in indicating whether differencing is needed. Recall that the ACF of an ARMA model should decay exponentially fast to zero. Slow decay in the sample ACF may be an indication that differencing is needed.

The next step is to look at the sample ACF and PACF of (possibly differenced) data. Using Table 3.1 as a guide, preliminary values of $p$ and $q$ are chosen. Recall that two models that are seemingly different can actually be very similar. With this in mind, we should not worry about being so precise at this stage of the model fitting. At this stage, a few preliminary values of $(p, d, q)$ should be at hand, and we can start fitting models.

[^16]

Fig. 3.13. Sample ACF of the GNP data. Lag is in terms of years.


Fig. 3.14. First difference of the U.S. GNP data.

## Example 3.30 Analysis of GNP Data

In this example, we consider the analysis of quarterly U.S. GNP from 1947(1) to 2002(3), $n=223$ observations. The data are real U.S. gross national product in billions of chained 1996 dollars and have been seasonally adjusted. The data were obtained from the Federal Reserve Bank of St. Louis (http://research.stlouisfed. org/). Figure 3.12 shows a plot of the data, say, $y_{t}$. Because strong trend hides any other effect, it is not clear from Figure 3.12 that the variance is increasing with time. For the purpose of demonstration, the sample ACF of the data is displayed in Figure 3.13. Figure 3.14 shows the first difference of the data, $\nabla y_{t}$, and now that the trend has been removed we are able to notice that the variability in the second half of the data is larger than in the first half of the data. Also, it appears as though a trend is still present after differencing. The growth rate, say, $x_{t}=\nabla \log \left(y_{t}\right)$, is plotted in Figure 3.15, and, appears to be a stable process. Moreover, we may interpret the values of $x_{t}$ as the percentage quarterly growth of U.S. GNP.

The sample ACF and PACF of the quarterly growth rate are plotted in Figure 3.16. Inspecting the sample ACF and PACF, we might feel that the ACF is cutting off at lag 2 and the PACF is tailing off. This would suggest the GNP growth rate follows an MA(2) process, or log GNP follows an $\operatorname{ARIMA}(0,1,2)$ model.


Fig. 3.15. U.S. GNP quarterly growth rate.

Rather than focus on one model, we will also suggest that it appears that the ACF is tailing off and the PACF is cutting off at lag 1. This suggests an $\operatorname{AR}(1)$ model for the growth rate, or ARIMA $(1,1,0)$ for $\log$ GNP. As a preliminary analysis, we will fit both models.

Using MLE to fit the MA(2) model for the growth rate, $x_{t}$, the estimated model is

$$
\begin{equation*}
x_{t}=.008_{(.001)}+.303_{(.065)} \widehat{w}_{t-1}+.204_{(.064)} \widehat{w}_{t-2}+\widehat{w}_{t} \tag{3.75}
\end{equation*}
$$

where $\widehat{\sigma}_{w}=.0094$ is based on 219 degrees of freedom. The values in parentheses are the corresponding estimated standard errors. All of the regression coefficients are significant, including the constant. We make a special note of this because, as a default, some computer packages do not fit a constant in a differenced model. That is, these packages assume, by default, that there is no drift. In this example, not including a constant leads to the wrong conclusions about the nature of the U.S. economy. Not including a constant assumes the average quarterly growth rate is zero, whereas the U.S. GNP average quarterly growth rate is about $1 \%$ (which can be seen easily in Figure 3.15). We leave it to the reader to investigate what happens when the constant is not included.

The estimated AR(1) model is

$$
\begin{equation*}
x_{t}=.008_{(.001)}(1-.347)+.347_{(.063)} x_{t-1}+\widehat{w}_{t} \tag{3.76}
\end{equation*}
$$

where $\widehat{\sigma}_{w}=.0095$ on 220 degrees of freedom; note that the constant in (3.76) is $.008(1-.347)=.005$.

We will discuss diagnostics next, but assuming both of these models fit well, how are we to reconcile the apparent differences of the estimated models (3.75) and (3.76)? In fact, the fitted models are nearly the same. To show this, consider an AR(1) model of the form in (3.76) without a constant term; that is,

$$
x_{t}=.35 x_{t-1}+w_{t},
$$

and write it in its causal form, $x_{t}=\sum_{j=0}^{\infty} \psi_{j} w_{t-j}$, where we recall $\psi_{j}=.35^{j}$. Thus, $\psi_{0}=1, \psi_{1}=.350, \psi_{2}=.123, \psi_{3}=.043, \psi_{4}=.015, \psi_{5}=.005, \psi_{6}=.002, \psi_{7}=$ $.001, \psi_{8}=0, \psi_{9}=0, \psi_{10}=0$, and so forth. Thus,


Fig. 3.16. Sample ACF and PACF of the GNP quarterly growth rate. Lag is in terms of years.

$$
x_{t} \approx .35 w_{t-1}+.12 w_{t-2}+w_{t}
$$

which is similar to the fitted MA(2) model in (3.75).
The analysis can be performed in R as follows. ${ }^{9}$

```
plot(gnp)
acf2(gnp, 50)
gnpgr = diff(log(gnp)) # growth rate
plot(gnpgr)
acf2(gnpgr, 24)
sarima(gnpgr, 1, 0, 0) # AR(1)
sarima(gnpgr, 0, 0, 2) # MA(2)
ARMAtoMA(ar=.35, ma=0, 10) # prints psi-weights
```

The next step in model fitting is diagnostics. This investigation includes the analysis of the residuals as well as model comparisons. Again, the first step involves a time plot of the innovations (or residuals), $x_{t}-\widehat{x}_{t}^{t-1}$, or of the standardized innovations

$$
\begin{equation*}
e_{t}=\left(x_{t}-\widehat{x}_{t}^{t-1}\right) / \sqrt{\widehat{P}_{t}^{t-1}} \tag{3.77}
\end{equation*}
$$

where $\widehat{x}_{t}^{t-1}$ is the one-step-ahead prediction of $x_{t}$ based on the fitted model and $\widehat{P}_{t}^{t-1}$ is the estimated one-step-ahead error variance. If the model fits well, the standardized residuals should behave as an iid sequence with mean zero and variance one. The time plot should be inspected for any obvious departures from this assumption. Unless the time series is Gaussian, it is not enough that the residuals are uncorrelated. It is possible in the non-Gaussian case to have an uncorrelated process for which contiguous values are highly dependent. As an example, we mention the family of GARCH models.

Investigation of marginal normality can be accomplished visually by looking at a histogram of the residuals. In addition to this, a normal probability plot or a Q-Q plot can help in identifying departures from normality. See Johnson and Wichern (1992, Chapter 4) for details of this test as well as additional tests for multivariate normality.

[^17]There are several tests of randomness, for example the runs test, that could be applied to the residuals. We could also inspect the sample autocorrelations of the residuals, say, $\widehat{\rho}_{e}(h)$, for any patterns or large values. Recall that, for a white noise sequence, the sample autocorrelations are approximately independently and normally distributed with zero means and variances $1 / n$. Hence, a good check on the correlation structure of the residuals is to plot $\widehat{\rho}_{e}(h)$ versus $h$ along with the error bounds of $\pm 2 / \sqrt{n}$. The residuals from a model fit, however, will not quite have the properties of a white noise sequence and the variance of $\widehat{\rho}_{e}(h)$ can be much less than $1 / n$. Details can be found in Box and Pierce (1970) and McLeod (1978). This part of the diagnostics can be viewed as a visual inspection of $\widehat{\rho}_{e}(h)$ with the main concern being the detection of obvious departures from the independence assumption.

In addition to plotting $\widehat{\rho}_{e}(h)$, we can perform a general test of whiteness that takes into consideration the magnitudes of $\widehat{\rho}_{e}(h)$ as a group. The Ljung-Box-Pierce Q -statistic given by

$$
\begin{equation*}
Q=n(n+2) \sum_{h=1}^{H} \frac{\widehat{\rho}_{e}^{2}(h)}{n-h} \tag{3.78}
\end{equation*}
$$

can be used to perform such a test. The value $H$ in (3.78) is chosen somewhat arbitrarily, typically, $H=20$. Under the null hypothesis of model adequacy, asymptotically $(n \rightarrow \infty), Q \sim \chi_{H-p-q}^{2}$. Thus, we would reject the null hypothesis at level $\alpha$ if the value of $Q$ exceeds the $(1-\alpha)$-quantile of the $\chi_{H-p-q}^{2}$ distribution. Details can be found in Box and Pierce (1970), Ljung and Box (1978), and Davies et al. (1977). The basic idea is that if $w_{t}$ is white noise, then by Property $1.2, n \widehat{\rho}_{w}^{2}(h)$, for $h=1, \ldots, H$, are asymptotically independent $\chi_{1}^{2}$ random variables. This means that $n \sum_{h=1}^{H} \widehat{\rho}_{w}^{2}(h)$ is approximately a $\chi_{H}^{2}$ random variable. Because the test involves the ACF of residuals from a model fit, there is a loss of $p+q$ degrees of freedom; the other values in (3.78) are used to adjust the statistic to better match the asymptotic chi-squared distribution.

## Example 3.31 Diagnostics for GNP Growth Rate Example

We will focus on the MA(2) fit from Example 3.30; the analysis of the $\operatorname{AR}(1)$ residuals is similar. Figure 3.17 displays a plot of the standardized residuals, the ACF of the residuals, a boxplot of the standardized residuals, and the p-values associated with the Q-statistic, (3.78), at lags $H=3$ through $H=20$ (with corresponding degrees of freedom $H-2$ ).

Inspection of the time plot of the standardized residuals in Figure 3.17 shows no obvious patterns. Notice that there are outliers, however, with a few values exceeding 3 standard deviations in magnitude. The ACF of the standardized residuals shows no apparent departure from the model assumptions, and the Q -statistic is never significant at the lags shown. The normal Q-Q plot of the residuals shows departure from normality at the tails due to the outliers that occurred primarily in the 1950s and the early 1980s.

The model appears to fit well except for the fact that a distribution with heavier tails than the normal distribution should be employed. The diagnostics shown in Figure 3.17 are a by-product of the sarima command from the previous example.


Fig. 3.17. Diagnostics of the residuals from $\mathrm{MA}(2)$ fit on GNP growth rate.

## Example 3.32 Diagnostics for the Glacial Varve Series

In Example 3.20, we fit an $\operatorname{ARIMA}(0,1,1)$ model to the logarithms of the glacial varve data and there appears to be a small amount of autocorrelation left in the residuals and the Q-tests are all significant; see Figure 3.18.

To adjust for this problem, we fit an $\operatorname{ARIMA}(1,1,1)$ to the logged varve data and obtained the estimates

$$
\widehat{\phi}=.23_{(.05)}, \widehat{\theta}=-.89_{(.03)}, \widehat{\sigma}_{w}^{2}=.23
$$

Hence the AR term is significant. The Q-statistic p-values for this model are also displayed in Figure 3.18, and it appears this model fits the data well.

As previously stated, the diagnostics are byproducts of the individual sarima runs. We note that we did not fit a constant in either model because there is no apparent drift in the differenced, logged varve series. This fact can be verified by noting the constant is not significant when the command no.constant=TRUE is removed in the code:


Fig. 3.18. $Q$-statistic $p$-values for the $\operatorname{ARIMA}(0,1,1)$ fit [top] and the $\operatorname{ARIMA}(1,1,1)$ fit [bottom] to the logged varve data.

```
sarima(log(varve), 0, 1, 1, no.constant=TRUE) # ARIMA(0,1,1)
sarima(log(varve), 1, 1, 1, no.constant=TRUE) # ARIMA(1,1,1)
```

In Example 3.30, we have two competing models, an $\operatorname{AR}(1)$ and an MA(2) on the GNP growth rate, that each appear to fit the data well. In addition, we might also consider that an $\operatorname{AR}(2)$ or an MA(3) might do better for forecasting. Perhaps combining both models, that is, fitting an $\operatorname{ARMA}(1,2)$ to the GNP growth rate, would be the best. As previously mentioned, we have to be concerned with overfitting the model; it is not always the case that more is better. Overfitting leads to less-precise estimators, and adding more parameters may fit the data better but may also lead to bad forecasts. This result is illustrated in the following example.

## Example 3.33 A Problem with Overfitting

Figure 3.19 shows the U.S. population by official census, every ten years from 1910 to 1990 , as points. If we use these nine observations to predict the future population, we can use an eight-degree polynomial so the fit to the nine observations is perfect. The model in this case is

$$
x_{t}=\beta_{0}+\beta_{1} t+\beta_{2} t^{2}+\cdots+\beta_{8} t^{8}+w_{t} .
$$

The fitted line, which is plotted in the figure, passes through the nine observations. The model predicts that the population of the United States will be close to zero in the year 2000, and will cross zero sometime in the year 2002!

The final step of model fitting is model choice or model selection. That is, we must decide which model we will retain for forecasting. The most popular techniques, AIC, AICc, and BIC, were described in $\S 2.2$ in the context of regression models.

## Example 3.34 Model Choice for the U.S. GNP Series

Returning to the analysis of the U.S. GNP data presented in Example 3.30 and Example 3.31, recall that two models, an AR(1) and an MA(2), fit the GNP growth


Fig. 3.19. A perfect fit and a terrible forecast.
rate well. To choose the final model, we compare the AIC, the AICc, and the BIC for both models. These values are a byproduct of the sarima runs displayed at the end of Example 3.30, but for convenience, we display them again here (recall the growth rate data are in gnpgr):

```
sarima(gnpgr, 1, 0, 0) # AR(1)
    $AIC: -8.294403 $AICc: -8.284898 $BIC: -9.263748
sarima(gnpgr, 0, 0, 2) # MA(2)
    $AIC: -8.297693 $AICc: -8.287854 $BIC: -9.251711
```

The AIC and AICc both prefer the MA(2) fit, whereas the BIC prefers the simpler AR(1) model. It is often the case that the BIC will select a model of smaller order than the AIC or AICc. In this case, it is reasonable to retain the AR(1) because pure autoregressive models are easier to work with.

### 3.9 Regression with Autocorrelated Errors

In $\S 2.2$, we covered the classical regression model with uncorrelated errors $w_{t}$. In this section, we discuss the modifications that might be considered when the errors are correlated. That is, consider the regression model

$$
\begin{equation*}
y_{t}=\beta_{1} z_{t 1}+\cdots+\beta_{r} z_{t r}+x_{t}=\sum_{j=1}^{r} \beta_{j} z_{t j}+x_{t} \tag{3.79}
\end{equation*}
$$

where $x_{t}$ is a process with some covariance function $\gamma_{x}(s, t)$. In ordinary least squares, the assumption is that $x_{t}$ is white Gaussian noise, in which case $\gamma_{x}(s, t)=0$ for $s \neq t$ and $\gamma_{x}(t, t)=\sigma^{2}$, independent of $t$. If this is not the case, then weighted least squares should be used. ${ }^{10}$

[^18]In the time series case, it is often possible to assume a stationary covariance structure for the error process $x_{t}$ that corresponds to a linear process and try to find an ARMA representation for $x_{t}$. For example, if we have a pure $\operatorname{AR}(p)$ error, then

$$
\phi(B) x_{t}=w_{t},
$$

and $\phi(B)=1-\phi_{1} B-\cdots-\phi_{p} B^{p}$ is the linear transformation that, when applied to the error process, produces the white noise $w_{t}$. Multiplying the regression equation through by the transformation $\phi(B)$ yields,

$$
\underbrace{\phi(B) y_{t}}_{y_{t}^{*}}=\sum_{j=1}^{r} \beta_{j} \underbrace{\phi(B) z_{t j}}_{z_{t j}^{*}}+\underbrace{\phi(B) x_{t}}_{w_{t}}
$$

and we are back to the linear regression model where the observations have been transformed so that $y_{t}^{*}=\phi(B) y_{t}$ is the dependent variable, $z_{t j}^{*}=\phi(B) z_{t j}$ for $j=$ $1, \ldots, r$, are the independent variables, but the $\beta \mathrm{s}$ are the same as in the original model. For example, if $p=1$, then $y_{t}^{*}=y_{t}-\phi y_{t-1}$ and $z_{t j}^{*}=z_{t j}-\phi z_{t-1, j}$.

In the AR case, we may set up the least squares problem as minimizing the error sum of squares

$$
S(\boldsymbol{\phi}, \boldsymbol{\beta})=\sum_{t=1}^{n} w_{t}^{2}=\sum_{t=1}^{n}\left[\phi(B) y_{t}-\sum_{j=1}^{r} \beta_{j} \phi(B) z_{t j}\right]^{2}
$$

with respect to all the parameters, $\boldsymbol{\phi}=\left\{\phi_{1}, \ldots, \phi_{p}\right\}$ and $\boldsymbol{\beta}=\left\{\beta_{1}, \ldots, \beta_{r}\right\}$. Of course, this is done using numerical methods.

If the error process is $\operatorname{ARMA}(p, q)$, i.e., $\phi(B) x_{t}=\theta(B) w_{t}$, then in the above discussion, we transform by $\pi(B) x_{t}=w_{t}$, where, recalling (3.16), $\pi(B)=\theta(B)^{-1} \phi(B)$. In this case the error sum of squares also depends on $\boldsymbol{\theta}=\left\{\theta_{1}, \ldots, \theta_{q}\right\}$ :

$$
S(\boldsymbol{\phi}, \boldsymbol{\theta}, \boldsymbol{\beta})=\sum_{t=1}^{n} w_{t}^{2}=\sum_{t=1}^{n}\left[\pi(B) y_{t}-\sum_{j=1}^{r} \beta_{j} \pi(B) z_{t j}\right]^{2}
$$

At this point, the main problem is that we do not typically know the behavior of the noise $x_{t}$ prior to the analysis. An easy way to tackle this problem was first presented in Cochrane and Orcutt (1949), and with the advent of cheap computing is modernized below:
(i) First, run an ordinary regression of $y_{t}$ on $z_{t 1}, \ldots, z_{t r}$ (acting as if the errors are uncorrelated). Retain the residuals, $\widehat{x}_{t}=y_{t}-\sum_{j=1}^{r} \widehat{\beta}_{j} z_{t j}$.
(ii) Identify ARMA model(s) for the residuals $\widehat{x}_{t}$.
(iii) Run weighted least squares (or MLE) on the regression model with autocorrelated errors using the model specified in step (ii).
(iv) Inspect the residuals $\widehat{w}_{t}$ for whiteness, and adjust the model if necessary.


Fig. 3.20. Sample ACF and PACF of the mortality residuals indicating an $\operatorname{AR}(2)$ process.

## Example 3.35 Mortality, Temperature and Pollution

We consider the analyses presented in Example 2.2, relating mean adjusted temperature $T_{t}$, and particulate levels $P_{t}$ to cardiovascular mortality $M_{t}$. We consider the regression model

$$
\begin{equation*}
M_{t}=\beta_{1}+\beta_{2} t+\beta_{3} T_{t}+\beta_{4} T_{t}^{2}+\beta_{5} P_{t}+x_{t}, \tag{3.80}
\end{equation*}
$$

where, for now, we assume that $x_{t}$ is white noise. The sample ACF and PACF of the residuals from the ordinary least squares fit of (3.80) are shown in Figure 3.20, and the results suggest an $\operatorname{AR}(2)$ model for the residuals.

Our next step is to fit the correlated error model (3.80), but where $x_{t}$ is $\operatorname{AR}(2)$,

$$
x_{t}=\phi_{1} x_{t-1}+\phi_{2} x_{t-2}+w_{t}
$$

and $w_{t}$ is white noise. The model can be fit using the arima function as follows (partial output shown).

```
trend = time(cmort); temp = tempr - mean(tempr); temp2 = temp^2
```

fit $=\operatorname{lm}$ (cmort $\sim$ trend + temp + temp2 + part, na.action=NULL)
acf2(resid(fit), 52) \# implies AR2
sarima(cmort, 2,0,0, xreg=cbind(trend,temp,temp2,part))

Coefficients:

| ar1 | ar2 | intercept | trend | temp | temp2 | part |
| :--- | ---: | ---: | ---: | ---: | ---: | ---: |
|  | 0.3848 | 0.4326 | 80.2116 | -1.5165 | -0.0190 | 0.0154 |
| 0.1545 |  |  |  |  |  |  |
| s.e. | 0.0436 | 0.0400 | 1.8072 | 0.4226 | 0.0495 | 0.0020 |
| sigma^2 | 0.0272 |  |  |  |  |  |
| stimated as | 26.01: loglikelihood | $=-1549.04$, | aic $=3114.07$ |  |  |  |

The residual analysis output from sarima (not shown) shows no obvious departure of the residuals from whiteness.

### 3.10 Seasonal ARIMA Models

In this section, we introduce several modifications made to the ARIMA model to account for seasonal and nonstationary behavior. Often, the dependence on the past tends to occur most strongly at multiples of some underlying seasonal lag $s$. For example, with monthly economic data, there is a strong yearly component occurring at lags that are multiples of $s=12$, because of the strong connections of all activity to the calendar year. Data taken quarterly will exhibit the yearly repetitive period at $s=4$ quarters. Natural phenomena such as temperature also have strong components corresponding to seasons. Hence, the natural variability of many physical, biological, and economic processes tends to match with seasonal fluctuations. Because of this, it is appropriate to introduce autoregressive and moving average polynomials that identify with the seasonal lags. The resulting pure seasonal autoregressive moving average model, say, $\operatorname{ARMA}(P, Q)_{s}$, then takes the form

$$
\begin{equation*}
\Phi_{P}\left(B^{s}\right) x_{t}=\Theta_{Q}\left(B^{s}\right) w_{t}, \tag{3.81}
\end{equation*}
$$

where the operators

$$
\begin{equation*}
\Phi_{P}\left(B^{s}\right)=1-\Phi_{1} B^{s}-\Phi_{2} B^{2 s}-\cdots-\Phi_{P} B^{P s} \tag{3.82}
\end{equation*}
$$

and

$$
\begin{equation*}
\Theta_{Q}\left(B^{s}\right)=1+\Theta_{1} B^{s}+\Theta_{2} B^{2 s}+\cdots+\Theta_{Q} B^{Q s} \tag{3.83}
\end{equation*}
$$

are the seasonal autoregressive operator and the seasonal moving average operator of orders $P$ and $Q$, respectively, with seasonal period $s$.

Analogous to the properties of nonseasonal ARMA models, the pure seasonal $\operatorname{ARMA}(P, Q)_{s}$ is causal only when the roots of $\Phi_{P}\left(z^{s}\right)$ lie outside the unit circle, and it is invertible only when the roots of $\Theta_{Q}\left(z^{S}\right)$ lie outside the unit circle.

## Example 3.36 A Seasonal AR Series

A first-order seasonal autoregressive series that might run over months could be written as

$$
\left(1-\Phi B^{12}\right) x_{t}=w_{t}
$$

or

$$
x_{t}=\Phi x_{t-12}+w_{t} .
$$

This model exhibits the series $x_{t}$ in terms of past lags at the multiple of the yearly seasonal period $s=12$ months. It is clear from the above form that estimation and forecasting for such a process involves only straightforward modifications of the unit lag case already treated. In particular, the causal condition requires $|\Phi|<1$.

We simulated 3 years of data from the model with $\Phi=.9$, and exhibit the theoretical ACF and PACF of the model. See Figure 3.21.
set.seed (666)
phi $=c(\operatorname{rep}(\theta, 11), .9)$
$\mathrm{sAR}=\operatorname{arima} . \operatorname{sim}(\operatorname{list}(\operatorname{order}=\mathrm{c}(12,0,0)$, ar=phi), $\mathrm{n}=37)$
sAR $=$ ts (sAR, freq=12)
layout (matrix (c(1,2, 1,3), nc=2))


Fig. 3.21. Data generated from a seasonal $(s=12) \operatorname{AR}(1)$, and the true ACF and PACF of the model $x_{t}=.9 x_{t-12}+w_{t}$.

```
par(mar=c (3,3,2,1), mgp=c(1.6,.6,0))
plot(sAR, axes=FALSE, main='seasonal AR(1)', xlab="year", type='c')
Months = c("J","F","M","A","M","J","J","A","S","0","N","D")
points(sAR, pch=Months, cex=1.25, font=4, col=1:4)
axis(1, 1:4); abline(v=1:4, lty=2, col='#cccccc')
axis(2); box()
ACF = ARMAacf(ar=phi, ma=0, 100)
PACF = ARMAacf(ar=phi, ma=0, 100, pacf=TRUE)
plot(ACF,type="h", xlab="lag", ylim=c(-.1,1)); abline(h=0)
plot(PACF, type="h", xlab="lag", ylim=c(-.1,1)); abline(h=0)
```

For the first-order seasonal $(s=12)$ MA model, $x_{t}=w_{t}+\Theta w_{t-12}$, it is easy to verify that

$$
\begin{aligned}
\gamma(0) & =\left(1+\Theta^{2}\right) \sigma^{2} \\
\gamma( \pm 12) & =\Theta \sigma^{2} \\
\gamma(h) & =0, \quad \text { otherwise. }
\end{aligned}
$$

Thus, the only nonzero correlation, aside from lag zero, is

$$
\rho( \pm 12)=\Theta /\left(1+\Theta^{2}\right)
$$

For the first-order seasonal $(s=12)$ AR model, using the techniques of the nonseasonal AR(1), we have

Table 3.3. Behavior of the ACF and PACF for Pure SARMA Models

|  | $\operatorname{AR}(P)_{s}$ | $\operatorname{MA}(Q)_{s}$ | $\operatorname{ARMA}(P, Q)_{s}$ |
| :---: | :---: | :---: | :---: |
| ACF* $^{*}$ | Tails off at lags $k s$, | Cuts off after | Tails off at |
|  | $k=1,2, \ldots$, | lag $Q s$ | lags $k s$ |
| PACF* $^{*}$ | Cuts off after | Tails off at lags $k s$ | Tails off at |
|  | lag $P s$ | $k=1,2, \ldots$, | lags $k s$ |

*The values at nonseasonal lags $h \neq k s$, for $k=1,2, \ldots$, are zero.

$$
\begin{aligned}
\gamma(0) & =\sigma^{2} /\left(1-\Phi^{2}\right) \\
\gamma( \pm 12 k) & =\sigma^{2} \Phi^{k} /\left(1-\Phi^{2}\right) \quad k=1,2, \ldots \\
\gamma(h) & =0, \quad \text { otherwise. }
\end{aligned}
$$

In this case, the only non-zero correlations are

$$
\rho( \pm 12 k)=\Phi^{k}, \quad k=0,1,2, \ldots
$$

These results can be verified using the general result that $\gamma(h)=\Phi \gamma(h-12)$, for $h \geq 1$. For example, when $h=1, \gamma(1)=\Phi \gamma(11)$, but when $h=11$, we have $\gamma(11)=\Phi \gamma(1)$, which implies that $\gamma(1)=\gamma(11)=0$. In addition to these results, the PACF have the analogous extensions from nonseasonal to seasonal models. These results are demonstrated in Figure 3.21.

As an initial diagnostic criterion, we can use the properties for the pure seasonal autoregressive and moving average series listed in Table 3.3. These properties may be considered as generalizations of the properties for nonseasonal models that were presented in Table 3.1.

In general, we can combine the seasonal and nonseasonal operators into a multiplicative seasonal autoregressive moving average model, denoted by $\operatorname{ARMA}(p, q) \times$ $(P, Q)_{s}$, and write

$$
\begin{equation*}
\Phi_{P}\left(B^{s}\right) \phi(B) x_{t}=\Theta_{Q}\left(B^{s}\right) \theta(B) w_{t} \tag{3.84}
\end{equation*}
$$

as the overall model. Although the diagnostic properties in Table 3.3 are not strictly true for the overall mixed model, the behavior of the ACF and PACF tends to show rough patterns of the indicated form. In fact, for mixed models, we tend to see a mixture of the facts listed in Table 3.1 and Table 3.3. In fitting such models, focusing on the seasonal autoregressive and moving average components first generally leads to more satisfactory results.

## Example 3.37 A Mixed Seasonal Model

Consider an $\operatorname{ARMA}(0,1) \times(1,0)_{12}$ model

$$
x_{t}=\Phi x_{t-12}+w_{t}+\theta w_{t-1},
$$

where $|\Phi|<1$ and $|\theta|<1$. Then, because $x_{t-12}, w_{t}$, and $w_{t-1}$ are uncorrelated, and $x_{t}$ is stationary, $\gamma(0)=\Phi^{2} \gamma(0)+\sigma_{w}^{2}+\theta^{2} \sigma_{w}^{2}$, or


Fig. 3.22. ACF and PACF of the mixed seasonal ARMA model $x_{t}=.8 x_{t-12}+w_{t}-.5 w_{t-1}$.

$$
\gamma(0)=\frac{1+\theta^{2}}{1-\Phi^{2}} \sigma_{w}^{2}
$$

In addition, multiplying the model by $x_{t-h}, h>0$, and taking expectations, we have $\gamma(1)=\Phi \gamma(11)+\theta \sigma_{w}^{2}$, and $\gamma(h)=\Phi \gamma(h-12)$, for $h \geq 2$. Thus, the ACF for this model is

$$
\begin{aligned}
\rho(12 h) & =\Phi^{h} \quad h=1,2, \ldots \\
\rho(12 h-1) & =\rho(12 h+1)=\frac{\theta}{1+\theta^{2}} \Phi^{h} \quad h=0,1,2, \ldots, \\
\rho(h) & =0, \quad \text { otherwise. }
\end{aligned}
$$

The ACF and PACF for this model, with $\Phi=.8$ and $\theta=-.5$, are shown in Figure 3.22. These type of correlation relationships, although idealized here, are typically seen with seasonal data.

To reproduce Figure 3.22 in R, use the following commands:

```
phi = c(rep (0,11),.8)
ACF = ARMAacf(ar=phi, ma=-.5, 50)[-1] # [-1] removes 0 lag
PACF = ARMAacf(ar=phi, ma=-.5, 50, pacf=TRUE)
par(mfrow=c(1,2))
plot(ACF, type="h", xlab="lag", ylim=c(-.4,.8)); abline(h=0)
plot(PACF, type="h", xlab="lag", ylim=c(-.4,.8)); abline(h=0)
```

Seasonal persistence occurs when the process is nearly periodic in the season. For example, with average monthly temperatures over the years, each January would be approximately the same, each February would be approximately the same, and so on. In this case, we might think of average monthly temperature $x_{t}$ as being modeled as

$$
x_{t}=S_{t}+w_{t}
$$

where $S_{t}$ is a seasonal component that varies a little from one year to the next, according to a random walk,

$$
S_{t}=S_{t-12}+v_{t}
$$

In this model, $w_{t}$ and $v_{t}$ are uncorrelated white noise processes. The tendency of data to follow this type of model will be exhibited in a sample ACF that is large and decays very slowly at lags $h=12 k$, for $k=1,2, \ldots$. If we subtract the effect of successive years from each other, we find that

$$
\left(1-B^{12}\right) x_{t}=x_{t}-x_{t-12}=v_{t}+w_{t}-w_{t-12}
$$

This model is a stationary $\mathrm{MA}(1)_{12}$, and its ACF will have a peak only at lag 12. In general, seasonal differencing can be indicated when the ACF decays slowly at multiples of some season $s$, but is negligible between the periods. Then, a seasonal difference of order D is defined as

$$
\begin{equation*}
\nabla_{s}^{D} x_{t}=\left(1-B^{s}\right)^{D} x_{t}, \tag{3.85}
\end{equation*}
$$

where $D=1,2, \ldots$, takes positive integer values. Typically, $D=1$ is sufficient to obtain seasonal stationarity. Incorporating these ideas into a general model leads to the following definition.

Definition 3.8 The multiplicative seasonal autoregressive integrated moving average model, or SARIMA model is given by

$$
\begin{equation*}
\Phi_{P}\left(B^{s}\right) \phi(B) \nabla_{s}^{D} \nabla^{d} x_{t}=\delta+\Theta_{Q}\left(B^{s}\right) \theta(B) w_{t}, \tag{3.86}
\end{equation*}
$$

where $w_{t}$ is the usual Gaussian white noise process. The general model is denoted as ARIMA $(p, d, q) \times(P, D, Q)_{s}$. The ordinary autoregressive and moving average components are represented by polynomials $\phi(B)$ and $\theta(B)$ of orders $p$ and $q$, respectively, and the seasonal autoregressive and moving average components by $\Phi_{P}\left(B^{s}\right)$ and $\Theta_{Q}\left(B^{s}\right)$ of orders $P$ and $Q$ and ordinary and seasonal difference components by $\nabla^{d}=(1-B)^{d}$ and $\nabla_{s}^{D}=\left(1-B^{s}\right)^{D}$.

## Example 3.38 An SARIMA Model

Consider the following model, which often provides a reasonable representation for seasonal, nonstationary, economic time series. We exhibit the equations for the model, denoted by $\operatorname{ARIMA}(0,1,1) \times(0,1,1)_{12}$ in the notation given above, where the seasonal fluctuations occur every 12 months. Then, with $\delta=0$, the model (3.86) becomes

$$
\nabla_{12} \nabla x_{t}=\Theta\left(B^{12}\right) \theta(B) w_{t}
$$

or

$$
\begin{equation*}
\left(1-B^{12}\right)(1-B) x_{t}=\left(1+\Theta B^{12}\right)(1+\theta B) w_{t} \tag{3.87}
\end{equation*}
$$

Expanding both sides of (3.87) leads to the representation

$$
\left(1-B-B^{12}+B^{13}\right) x_{t}=\left(1+\theta B+\Theta B^{12}+\Theta \theta B^{13}\right) w_{t}
$$

or in difference equation form


Fig. 3.23. $R$ data set AirPassengers, which are the monthly totals of international airline passengers x , and the transformed data: $1 \mathrm{x}=\log x_{t}, \mathrm{dlx}=\nabla \log x_{t}$, and ddlx $=\nabla_{12} \nabla \log x_{t}$.

$$
x_{t}=x_{t-1}+x_{t-12}-x_{t-13}+w_{t}+\theta w_{t-1}+\Theta w_{t-12}+\Theta \theta w_{t-13} .
$$

Note that the multiplicative nature of the model implies that the coefficient of $w_{t-13}$ is the product of the coefficients of $w_{t-1}$ and $w_{t-12}$ rather than a free parameter. The multiplicative model assumption seems to work well with many seasonal time series data sets while reducing the number of parameters that must be estimated.

Selecting the appropriate model for a given set of data from all of those represented by the general form (3.86) is a daunting task, and we usually think first in terms of finding difference operators that produce a roughly stationary series and then in terms of finding a set of simple autoregressive moving average or multiplicative seasonal ARMA to fit the resulting residual series. Differencing operations are applied first, and then the residuals are constructed from a series of reduced length. Next, the ACF and the PACF of these residuals are evaluated. Peaks that appear in these functions can often be eliminated by fitting an autoregressive or moving average component in accordance with the general properties of Table 3.1 and Table 3.3. In considering whether the model is satisfactory, the diagnostic techniques discussed in $\S 3.8$ still apply.

## Example 3.39 Air Passengers

We consider the R data set AirPassengers, which are the monthly totals of international airline passengers, 1949 to 1960, taken from Box \& Jenkins (1970). Various plots of the data and transformed data are shown in Figure 3.23 and were obtained as follows:

```
x = AirPassengers
lx = log(x); dlx = diff(lx); ddlx = diff(dlx, 12)
plot.ts(cbind(x,lx,dlx,ddlx), main="")
# of interest for showing seasonal RW (not shown here):
dev.new()
par(mfrow=c (2,1), mar=c (3,3,1,1), mgp=c (1.6,.6,0))
monthplot(dlx)
monthplot(ddlx)
```

Note that x is the original series, which shows trend plus increasing variance. The logged data are in 1 x , and the transformation stabilizes the variance. The logged data are then differenced to remove trend, and are stored in dlx. It is clear the there is still persistence in the seasons (i.e., $\mathrm{dl} \mathrm{x}_{t} \approx \mathrm{dl} \mathrm{x}_{t-12}$ ), so that a twelfth-order difference is applied and stored in ddlx. The transformed data appears to be stationary and we are now ready to fit a model.

The sample ACF and PACF of ddlx $\left(\nabla_{12} \nabla \log x_{t}\right)$ are shown in Figure 3.24. The R code is:

```
acf2(ddlx,50)
```

Between Seasons: It appears that at the seasons, the ACF is cutting off a lag $1 s$ $(s=12)$, whereas the PACF is tailing off at lags $1 s, 2 s, 3 s, 4 s, \ldots$. These results implies an $\operatorname{SMA}(1), P=0, Q=1$, in the season ( $s=12$ ).
Within Seasons: Inspecting the sample ACF and PACF at the lower lags, it appears as though both are tailing off. This suggests an $\operatorname{ARMA}(1,1)$ within the seasons, $p=q=1$.

Thus, we first try an $\operatorname{ARIMA}(1,1,1) \times(0,1,1)_{12}$ on the logged data:
$\operatorname{sarima}(1 \mathrm{x}, 1,1,1,0,1,1,12)$
Coefficients:

| ar1 | ma1 | sma1 |  |
| :--- | ---: | ---: | ---: |
|  | 0.1960 | -0.5784 | -0.5643 |
| s.e. | 0.2475 | 0.2132 | 0.0747 |
| sigma^2 estimated as 0.001341 |  |  |  |
| \$AIC | -5.5726 | \$AICc | -5.556713 | \$BIC -6.510729

However, the AR parameter is not significant, so we should try dropping one parameter from the within seasons part. In this case, we try both an $\operatorname{ARIMA}(0,1,1) \times$ $(0,1,1)_{12}$ and an $\operatorname{ARIMA}(1,1,0) \times(0,1,1)_{12}$ model:
sarima(lx, $0,1,1,0,1,1,12$ )
Coefficients:

$$
\begin{array}{rrr} 
& \mathrm{ma1} & \mathrm{sma1} \\
& -0.4018 & -0.5569 \\
\text { s.e. } & 0.0896 & 0.0731
\end{array}
$$

sigma^2 estimated as 0.001348
\$AIC -5.58133 \$AICc -5.56625 \$BIC -6.540082


Fig. 3.24. Sample ACF and PACF of ddlx $\left(\nabla_{12} \nabla \log x_{t}\right)$.

```
sarima(lx, 1,1,0, 0,1,1,12)
    Coefficients:
            ar1 sma1
        -0.3395 -0.5619
s.e. 0.0822 0.0748
    sigma^2 estimated as 0.001367
    $AIC -5.567081 $AICC -5.552002 $BIC -6.525834
```

All information criteria prefer the $\operatorname{ARIMA}(0,1,1) \times(0,1,1)_{12}$ model, which is the model displayed in (3.87). The residual diagnostics are shown in Figure 3.25, and except for one or two outliers, the model seems to fit well.

Finally, we forecast the logged data out twelve months, and the results are shown in Figure 3.26.
sarima. for (lx, $12,0,1,1,0,1,1,12$ )

## Example 3.40 Regression with Lagged Variables (cont)

In Example 2.8 we fit the model

$$
R_{t}=\beta_{1}+\beta_{2} S_{t-6}+\beta_{3} D_{t-6}+\beta_{4} D_{t-6} S_{t-6}+w_{t},
$$

where $R_{t}$ is Recruitment, $S_{t}$ is SOI, and $D_{t}$ is a dummy variable that is 0 if $S_{t}<0$ and 1 otherwise. However, residual analysis indicates that the residuals are not white noise. The sample (P)ACF of the residuals indicates that an AR(2) model might be appropriate, which is similar to the results of Example 3.35. Additional analysis suggests that at seasonal model might be needed, and after some investigation, we settled on an $\operatorname{SARIMA}(2,0,0) \times(2,0,0)_{12}$ model for the residuals. Both AIC and


Fig. 3.25. Residual analysis for the $\operatorname{ARIMA}(0,1,1) \times(0,1,1)_{12}$ fit to the logged air passengers data set.

BIC confirm that this model is the the best. We only display the results of the final model for this example, but the code to carry out the complete analysis is listed below.

```
dummy=ifelse(soi<0,0,1)
fish=ts.intersect(rec, soiL6=lag(soi,-6),dL6=lag(dummy,-6),dframe=TRUE)
summary(fit<-lm(rec~soiL6*dL6, data=fish, na.action=NULL))
attach(fish)
plot(resid(fit))
acf2(resid(fit)) # indicates AR(2)
(fit1 = sarima(rec,2,0,0, xreg=cbind(soiL6,dL6,I(soiL6*dL6))))
acf2(resid(fit1)) # indicates seasonal AR
# try seasonal AR order 1 and then 2, which appears to be best
(fit2 = sarima(rec, 2,0,0,1,0,0,12,xreg=cbind(soiL6,dL6,I(soiL6*dL6))))
(fit3 = sarima(rec,2,0,0,2,0,0,12,xreg=cbind(soiL6,dL6,I(soiL6*dL6))))
Coefficients: (rounded)
ar1 ar2 sar1 sar2 intrcpt soiL6 dL6 I(soiL6*dL6)
    1.349 -0.444 0.114 0.15 64 8.4 -2.25 
s.e. 0.044 0.045 0.047 0.05 < 6 2.2 0.95 2.9
sigma^2 estimated as 83.75: log likelihood = -1625.53, aic = 3269.07
$AIC $AICc $BIC
[1] 5.463671 [1] 5.469067 [1] 4.537095
```



Fig. 3.26. Twelve month forecast using the $\operatorname{ARIMA}(0,1,1) \times(0,1,1)_{12}$ model on the logged air passenger data set.

## Problems

3.1 For an MA(1), $x_{t}=w_{t}+\theta w_{t-1}$, show that $\left|\rho_{x}(1)\right| \leq 1 / 2$ for any number $\theta$. For which values of $\theta$ does $\rho_{x}(1)$ attain its maximum and minimum?
3.2 Let $\left\{w_{t} ; t=0,1, \ldots\right\}$ be a white noise process with variance $\sigma_{w}^{2}$ and let $|\phi|<1$ be a constant. Consider the process $x_{0}=w_{0}$, and

$$
x_{t}=\phi x_{t-1}+w_{t}, \quad t=1,2, \ldots .
$$

We might use this method to simulate an $\operatorname{AR}(1)$ process from simulated white noise.
(a) Show that $x_{t}=\sum_{j=0}^{t} \phi^{j} w_{t-j}$ for any $t=0,1, \ldots$
(b) Find the $E\left(x_{t}\right)$.
(c) Show that, for $t=0,1, \ldots,{ }^{11}$

$$
\operatorname{var}\left(x_{t}\right)=\frac{\sigma_{w}^{2}}{1-\phi^{2}}\left(1-\phi^{2(t+1)}\right)
$$

(d) Show that, for $h \geq 0,12$

$$
\operatorname{cov}\left(x_{t+h}, x_{t}\right)=\phi^{h} \operatorname{var}\left(x_{t}\right)
$$

(e) Is $x_{t}$ stationary?
(f) Argue that, as $t \rightarrow \infty$, the process becomes stationary, so in a sense, $x_{t}$ is "asymptotically stationary."
(g) Comment on how you could use these results to simulate $n$ observations of a stationary Gaussian $\operatorname{AR}(1)$ model from simulated iid $\mathrm{N}(0,1)$ values.
(h) Now suppose $x_{0}=w_{0} / \sqrt{1-\phi^{2}}$. Is this process stationary? Hint: Show $\operatorname{var}\left(x_{t}\right)$ is constant.
${ }^{11} \sum_{j=0}^{k} a^{j}=\left(1-a^{k+1}\right) /(1-a)$ for $|a| \neq 1$
${ }^{12}$ Use footnote 1 to write $x_{t+h}$ in terms of $x_{t}$ and other stuff.
3.3 Using Example 3.6 as a guide, identify the following models as $\operatorname{ARMA}(p, q)$ models (watch out for parameter redundancy), and determine whether they are causal and/or invertible. If the model is causal, use R to find the first $10 \psi$-weights, and if the model is invertible, use R to find the first $10 \pi$-weights.
(a) $x_{t}=.80 x_{t-1}-.15 x_{t-2}+w_{t}-.30 w_{t-1}$.
(b) $x_{t}=x_{t-1}-.50 x_{t-2}+w_{t}-w_{t-1}$.
3.4 For the $\mathrm{AR}(2)$ model given by $x_{t}=-.9 x_{t-2}+w_{t}$, follow the R code in Example 3.11 to find the roots of the autoregressive polynomial, find the pseudo period of the process, and then plot the theoretical ACF, $\rho(h)$.
3.5 (a) Compare the theoretical ACF and PACF of an ARMA(1, 1), an ARMA $(1,0)$, and an $\operatorname{ARMA}(0,1)$ series by plotting the ACFs and PACFs of the three series for $\phi=.6, \theta=.9$. Comment on the capability of the ACF and PACF to determine the order of the models. Hint: See the code for Example 3.13.
(b) Use arima.sim to generate $n=100$ observations from each of the three models discussed in (a). Compute the sample ACFs and PACFs for each model and compare it to the theoretical values. How do the results compare with the general results given in Table 3.1?
(c) Repeat (b) but with $n=500$. Comment.
3.6 Let $c_{t}$ be the cardiovascular mortality series (cmort) discussed in Chapter 2, Example 2.2 and let $x_{t}=\nabla c_{t}$ be the differenced data.
(a) Plot $x_{t}$ and compare it to the actual data plotted in Figure 2.2. Why does differencing seem reasonable in this case?
(b) Calculate and plot the sample ACF and PACF of $x_{t}$ and using Table 3.1, argue that an $\mathrm{AR}(1)$ is appropriate for $x_{t}$.
(c) Fit an $\operatorname{AR}(1)$ to $x_{t}$ using maximum likelihood (basically unconditional least squares) as in Section 3.6. The easiest way to do this is to use sarima from astsa. Comment on the significance of the regression parameter estimates of the model. What is the estimate of the white noise variance?
(d) Examine the residuals and comment on whether or not you think the residuals are white.
(e) Assuming the fitted model is the true model, find the forecasts over a four-week horizon, $x_{n+m}^{n}$, for $m=1,2,3,4$, and the corresponding $95 \%$ prediction intervals; $n=508$ here. The easiest way to do this is to use sarima.for from astsa.
(f) Show how the values obtained in part (e) were calculated.
(g) What is the one-step-ahead forecast of the actual value of cardiovascular mortality; i.e., what is $c_{n+1}^{n}$ ?
3.7 For an $\operatorname{AR}(1)$ model, determine the general form of the $m$-step-ahead forecast $x_{n+m}^{n}$ and show

$$
E\left[\left(x_{n+m}-x_{n+m}^{n}\right)^{2}\right]=\sigma_{w}^{2} \frac{1-\phi^{2 m}}{1-\phi^{2}}
$$

3.8 Repeat the following numerical exercise five times. Generate $n=100$ iid $\mathrm{N}(0,1)$ observations. Fit an $\operatorname{ARMA}(1,1)$ model to the data. Compare the parameter estimates in each case and explain the results.
3.9 Generate 10 realizations of length $n=200$ each of an $\operatorname{ARMA}(1,1)$ process with $\phi=.9, \theta=.5$ and $\sigma^{2}=1$. Find the MLEs of the three parameters in each case and compare the estimators to the true values.
3.10 Using Example 3.19 as your guide, find the Gauss-Newton procedure for estimating the autoregressive parameter, $\phi$, from the $\operatorname{AR}(1)$ model, $x_{t}=\phi x_{t-1}+w_{t}$, given data $x_{1}, \ldots, x_{n}$. Does this procedure produce the unconditional or the conditional estimator? Hint: Write the model as $w_{t}(\phi)=x_{t}-\phi x_{t-1}$; your solution should work out to be a non-recursive procedure.
3.11 Verify that the $\operatorname{IMA}(1,1)$ model given in (3.72) can be inverted and written as (3.73).
3.12 For the logarithm of the glacial varve data, say, $x_{t}$, presented in Example 3.20, use the first 100 observations and calculate the EWMA, $x_{n+1}^{n}$, discussed in Example 3.29, for $n=1, \ldots, 100$, using $\lambda=.25, .50$, and .75 , and plot the EWMAs and the data superimposed on each other. Comment on the results.
3.13 Crude oil prices in dollars per barrel are in oil; see Appendix R for more details. Fit an $\operatorname{ARIMA}(p, d, q)$ model to the growth rate performing all necessary diagnostics. Comment.
3.14 Fit an $\operatorname{ARIMA}(p, d, q)$ model to the global temperature data gtemp2 in astsa performing all of the necessary diagnostics. After deciding on an appropriate model, forecast (with limits) the next 10 years. Comment.
3.15 One of the series collected along with particulates, temperature, and mortality described in Example 2.2 is the sulfur dioxide series, so2. Fit an $\operatorname{ARIMA}(p, d, q)$ model to the data, performing all of the necessary diagnostics. After deciding on an appropriate model, forecast the data into the future four time periods ahead (about one month) and calculate $95 \%$ prediction intervals for each of the four forecasts. Comment.
3.16 Let $S_{t}$ represent the monthly sales data in sales ( $n=150$ ), and let $L_{t}$ be the leading indicator in lead.
(a) Fit an ARIMA model to $S_{t}$, the monthly sales data. Discuss your model fitting in a step-by-step fashion, presenting your (A) initial examination of the data, (B) transformations, if necessary, (C) initial identification of the dependence orders and degree of differencing, (D) parameter estimation, (E) residual diagnostics and model choice.
(b) Use the CCF and lag plots between $\nabla S_{t}$ and $\nabla L_{t}$ to argue that a regression of $\nabla S_{t}$ on $\nabla L_{t-3}$ is reasonable. [Note: In lag2.plot(), the first named series is the one that gets lagged.]
(c) Fit the regression model $\nabla S_{t}=\beta_{0}+\beta_{1} \nabla L_{t-3}+x_{t}$, where $x_{t}$ is an ARMA process (explain how you decided on your model for $x_{t}$ ). Discuss your results. $R$ help: If you have to work with various transformations of series in x and y , first align the data:
$\operatorname{dog}=$ ts.intersect ( $\operatorname{lag}(x,-11), \operatorname{diff}(y, 97)$ )
xnew $=\operatorname{dog}[, 1]$ \# dog has 2 columns, the first is lag(x,-11) ...
ynew $=\operatorname{dog}[, 2]$ \# ... and the second column is $\operatorname{diff}(y, 97)$
plot(dog) \# now you can manipulate xnew and ynew simultaneously lag2.plot(xnew, ynew, 5)
3.17 Redo Problem 2.2 without assuming the error term is white noise.
3.18 Plot the theoretical ACF of the seasonal $\operatorname{ARIMA}(0,1) \times(1,0)_{12}$ model with $\Phi=.8$ and $\theta=.5$ out to lag 50 .
3.19 Fit a seasonal ARIMA model of your choice to the unemployment data, unemp. Use the estimated model to forecast the next 12 months.
3.20 Fit a seasonal ARIMA model of your choice to the U.S. Live Birth Series, birth. Use the estimated model to forecast the next 12 months.
3.21 Fit an appropriate seasonal ARIMA model to the log-transformed Johnson and Johnson earnings series ( jj ) of Example 1.1. Use the estimated model to forecast the next 4 quarters.

## Chapter 4

## Spectral Analysis and Filtering

### 4.1 Introduction

The cyclic behavior of data is the focus of this chapter. For example, in the Johnson \& Johnson data set in Figure 1.1, the predominant frequency of oscillation is one cycle per year (4 quarters), or .25 cycles per observation. The predominant frequency in the SOI and fish populations series in Figure 1.5 is also one cycle per year, but this corresponds to 1 cycle every 12 months, or .083 cycles per observation. For simplicity, we measure frequency, $\omega$, at cycles per time point and discuss the implications of certain frequencies in terms of the problem context. Of descriptive interest is the period of a time series, defined as the number of points in a cycle, i.e., $1 / \omega$. Hence, the predominant period of the Johnson \& Johnson series is $1 / .25$ or 4 quarters per cycle, whereas the predominant period of the SOI series is 12 months per cycle. As stated in the Preface, complex numbers (a pdf) may be helpful for this chapter.

### 4.2 Periodicity and Cyclical Behavior

We have already encountered the notion of periodicity in numerous examples in Chapters 1, 2 and 3. The general notion of periodicity can be made more precise by introducing some terminology. In order to define the rate at which a series oscillates, we first define a cycle as one complete period of a sine or cosine function defined over a unit time interval. As in (1.5), we consider the periodic process

$$
\begin{equation*}
x_{t}=A \cos (2 \pi \omega t+\phi) \tag{4.1}
\end{equation*}
$$

for $t=0, \pm 1, \pm 2, \ldots$, where $\omega$ is a frequency index, defined in cycles per unit time with $A$ determining the height or amplitude of the function and $\phi$, called the phase, determining the start point of the cosine function. We can introduce random variation in this time series by allowing the amplitude and phase to vary randomly.

As discussed in Example 2.9, for purposes of data analysis, it is easier to use a trigonometric identity ${ }^{1}$ and write (4.1) as

[^19]\[

$$
\begin{equation*}
x_{t}=U_{1} \cos (2 \pi \omega t)+U_{2} \sin (2 \pi \omega t), \tag{4.2}
\end{equation*}
$$

\]

where $U_{1}=A \cos \phi$ and $U_{2}=-A \sin \phi$ are often taken to be normally distributed random variables. In this case, the amplitude is $A=\sqrt{ }\left(U_{1}^{2}+U_{2}^{2}\right)$ and the phase is $\phi=\tan ^{-1}\left(-U_{2} / U_{1}\right)$. From these facts we can show that if, and only if, in (4.1), $A$ and $\phi$ are independent random variables, where $A^{2}$ is chi-squared with 2 degrees of freedom, and $\phi$ is uniformly distributed on $(-\pi, \pi)$, then $U_{1}$ and $U_{2}$ are independent, standard normal random variables.

If we assume that $U_{1}$ and $U_{2}$ are iid random variables with mean 0 and variance $\sigma^{2}$, then $x_{t}$ in (4.2) is stationary with mean $E\left(x_{t}\right)=0$ and, writing $c_{t}=\cos (2 \pi \omega t)$ and $s_{t}=\sin (2 \pi \omega t)$, autocovariance function

$$
\begin{align*}
\gamma(h)=\operatorname{cov}\left(x_{t+h}, x_{t}\right)= & \operatorname{cov}\left(U_{1} c_{t+h}+U_{2} s_{t+h}, U_{1} c_{t}+U_{2} s_{t}\right) \\
& =\operatorname{cov}\left(U_{1} c_{t+h}, U_{1} c_{t}\right)+\operatorname{cov}\left(U_{1} c_{t+h}, U_{2} s_{t}\right) \\
& +\operatorname{cov}\left(U_{2} s_{t+h}, U_{1} c_{t}\right)+\operatorname{cov}\left(U_{2} s_{t+h}, U_{2} s_{t}\right)  \tag{4.3}\\
& =\sigma^{2} c_{t+h} c_{t}+0+0+\sigma^{2} s_{t+h} s_{t}=\sigma^{2} \cos (2 \pi \omega h),
\end{align*}
$$

using footnote 1 and noting that $\operatorname{cov}\left(U_{1}, U_{2}\right)=0$.
The random process in (4.2) is function of its frequency, $\omega$. For $\omega=1$, the series makes one cycle per time unit; for $\omega=.50$, the series makes a cycle every two time units; for $\omega=.25$, every four units, and so on. In general, for data that occur at discrete time points, we will need at least two points to determine a cycle, so the highest frequency of interest is .5 cycles per point. This frequency is called the folding frequency and defines the highest frequency that can be seen in discrete sampling. Higher frequencies sampled this way will appear at lower frequencies, called aliases; an example is the way a camera samples a rotating wheel on a moving automobile in a movie, in which the wheel appears to be rotating at a different rate. For example, movies are recorded at 24 frames per second. If the camera is filming a wheel that is rotating at the rate of 24 cycles per second (or 24 Hertz), the wheel will appear to stand still.

Consider a generalization of (4.2) that allows mixtures of periodic series with multiple frequencies and amplitudes,

$$
\begin{equation*}
x_{t}=\sum_{k=1}^{q}\left[U_{k 1} \cos \left(2 \pi \omega_{k} t\right)+U_{k 2} \sin \left(2 \pi \omega_{k} t\right)\right], \tag{4.4}
\end{equation*}
$$

where $U_{k 1}, U_{k 2}$, for $k=1,2, \ldots, q$, are independent zero-mean random variables with variances $\sigma_{k}^{2}$, and the $\omega_{k}$ are distinct frequencies. Notice that (4.4) exhibits the process as a sum of independent components, with variance $\sigma_{k}^{2}$ for frequency $\omega_{k}$. As in (4.3), it is easy to show (Problem 4.2) that the autocovariance function of the process is

$$
\begin{equation*}
\gamma(h)=\sum_{k=1}^{q} \sigma_{k}^{2} \cos \left(2 \pi \omega_{k} h\right), \tag{4.5}
\end{equation*}
$$

and we note the autocovariance function is the sum of periodic components with weights proportional to the variances $\sigma_{k}^{2}$. Hence, $x_{t}$ is a mean-zero stationary processes with variance


Fig. 4.1. Periodic components and their sum as described in Example 4.1.

$$
\begin{equation*}
\gamma(0)=\operatorname{var}\left(x_{t}\right)=\sum_{k=1}^{q} \sigma_{k}^{2}, \tag{4.6}
\end{equation*}
$$

which exhibits the overall variance as a sum of variances of each of the component parts.

## Example 4.1 A Periodic Series

Figure 4.1 shows an example of the mixture (4.4) with $q=3$ constructed in the following way. First, for $t=1, \ldots, 100$, we generated three series

$$
\begin{aligned}
& x_{t 1}=2 \cos (2 \pi t 6 / 100)+3 \sin (2 \pi t 6 / 100) \\
& x_{t 2}=4 \cos (2 \pi t 10 / 100)+5 \sin (2 \pi t 10 / 100) \\
& x_{t 3}=6 \cos (2 \pi t 40 / 100)+7 \sin (2 \pi t 40 / 100)
\end{aligned}
$$

These three series are displayed in Figure 4.1 along with the corresponding frequencies and squared amplitudes. For example, the squared amplitude of $x_{t 1}$ is $A^{2}=2^{2}+3^{2}=13$. Hence, the maximum and minimum values that $x_{t 1}$ will attain are $\pm \sqrt{13}= \pm 3.61$.

Finally, we constructed

$$
x_{t}=x_{t 1}+x_{t 2}+x_{t 3}
$$

and this series is also displayed in Figure 4.1. We note that $x_{t}$ appears to behave as some of the periodic series we saw in Chapters 1 and 2. The systematic sorting out of the essential frequency components in a time series, including their relative contributions, constitutes one of the main objectives of spectral analysis.

The R code to reproduce Figure 4.1 is

```
x1 = 2*cos(2*pi*1:100*6/100) + 3*sin(2*pi*1:100*6/100)
x2 = 4*cos(2*pi*1:100*10/100) + 5*sin(2*pi*1:100*10/100)
x3 = 6*cos(2*pi*1:100*40/100) + 7*sin(2*pi*1:100*40/100)
x = x1 + x2 + x3
par(mfrow=c(2,2))
plot.ts(x1, ylim=c(-10,10), main=expression(omega==6/100~~~A^2==13))
plot.ts(x2, ylim=c(-10,10), main=expression(omega==10/100~~~A^2==41))
plot.ts(x3, ylim=c(-10,10), main=expression(omega==40/100~~~A^2==85))
plot.ts(x, ylim=c(-16,16), main="sum")
```

The model given in (4.4), along with its autocovariance given (4.5), is a population construct. If the model is correct, our next step would be to estimate the variances $\sigma_{k}^{2}$ and frequencies $\omega_{k}$ that form the model (4.4). In the next example, we consider the problem of estimation of these quantities.

## Example 4.2 Estimation and the Periodogram

For any time series sample $x_{1}, \ldots, x_{n}$, where $n$ is odd, we may write, exactly

$$
\begin{equation*}
x_{t}=a_{0}+\sum_{j=1}^{(n-1) / 2}\left[a_{j} \cos (2 \pi t j / n)+b_{j} \sin (2 \pi t j / n)\right] \tag{4.7}
\end{equation*}
$$

for $t=1, \ldots, n$ and suitably chosen coefficients. If $n$ is even, the representation (4.7) can be modified by summing to ( $n / 2-1$ ) and adding an additional component given by $a_{n / 2} \cos \left(2 \pi t \frac{1}{2}\right)=a_{n / 2}(-1)^{t}$. The crucial point here is that (4.7) is exact for any sample. Hence (4.4) may be thought of as an approximation to (4.7), the idea being that many of the coefficients in (4.7) may be close to zero.

Using the regression results from Chapter 2, the coefficients $a_{j}$ and $b_{j}$ are of the form $\sum_{t=1}^{n} x_{t} z_{t j} / \sum_{t=1}^{n} z_{t j}^{2}$, where $z_{t j}$ is either $\cos (2 \pi t j / n)$ or $\sin (2 \pi t j / n)$. Using Problem 4.20, $\sum_{t=1}^{n} z_{t j}^{2}=n / 2$ when $j / n \neq 0,1 / 2$, so the regression coefficients in (4.7) can be written as

$$
a_{j}=\frac{2}{n} \sum_{t=1}^{n} x_{t} \cos (2 \pi t j / n) \quad \text { and } \quad b_{j}=\frac{2}{n} \sum_{t=1}^{n} x_{t} \sin (2 \pi t j / n)
$$

We then define the scaled periodogram to be

$$
\begin{equation*}
P(j / n)=a_{j}^{2}+b_{j}^{2} \tag{4.8}
\end{equation*}
$$

because it indicates which frequency components in (4.7) are large in magnitude and which components are small. The scaled periodogram is the estimate of $\sigma_{j}^{2}$ corresponding to the sinusoid oscillating at a frequency of $\omega_{j}=j / n$, or $j$ cycles in $n$ time points. These particular frequencies are called the Fourier or fundamental frequencies. Large values of $P(j / n)$ indicate which frequencies $\omega_{j}=j / n$ are predominant in the series, whereas small values of $P(j / n)$ may be associated with noise.

It is not necessary to run a large regression to obtain the values of $a_{j}$ and $b_{j}$ because they can be computed quickly if $n$ is a highly composite integer. Although


Fig. 4.2. Periodogram of the data generated in Example 4.1.
we will discuss it in more detail in $\S 4.4$, the discrete Fourier transform (DFT) is a complex-valued weighted average of the data given by ${ }^{2}$

$$
\begin{align*}
d(j / n) & =n^{-1 / 2} \sum_{t=1}^{n} x_{t} \exp (-2 \pi i t j / n)  \tag{4.9}\\
& =n^{-1 / 2}\left(\sum_{t=1}^{n} x_{t} \cos (2 \pi t j / n)-i \sum_{t=1}^{n} x_{t} \sin (2 \pi t j / n)\right)
\end{align*}
$$

for $j=0,1, \ldots, n-1$, where the frequencies $j / n$ are called the Fourier or fundamental frequencies. Because of a large number of redundancies in the calculation, (4.9) may be computed quickly using the fast Fourier transform (FFT). Note that

$$
\begin{equation*}
|d(j / n)|^{2}=\frac{1}{n}\left(\sum_{t=1}^{n} x_{t} \cos (2 \pi t j / n)\right)^{2}+\frac{1}{n}\left(\sum_{t=1}^{n} x_{t} \sin (2 \pi t j / n)\right)^{2} \tag{4.10}
\end{equation*}
$$

and it is this quantity that is called the periodogram. We may calculate the scaled periodogram, (4.8), using the periodogram as

$$
\begin{equation*}
P(j / n)=\frac{4}{n}|d(j / n)|^{2} . \tag{4.11}
\end{equation*}
$$

The scaled periodogram of the data, $x_{t}$, simulated in Example 4.1 is shown in Figure 4.2 , and it clearly identifies the three components $x_{t 1}, x_{t 2}$, and $x_{t 3}$ of $x_{t}$. Note that

$$
P(j / n)=P(1-j / n), \quad j=0,1, \ldots, n-1,
$$

so there is a mirroring effect at the folding frequency of $1 / 2$; consequently, the periodogram is typically not plotted for frequencies higher than the folding frequency. In addition, note that the heights of the scaled periodogram shown in the figure are

[^20]$$
P\left(\frac{6}{100}\right)=P\left(\frac{94}{100}\right)=13, \quad P\left(\frac{10}{100}\right)=P\left(\frac{90}{100}\right)=41, \quad P\left(\frac{40}{100}\right)=P\left(\frac{60}{100}\right)=85,
$$
and $P(j / n)=0$ otherwise. These are exactly the values of the squared amplitudes of the components generated in Example 4.1. This outcome suggests that the periodogram may provide some insight into the variance components, (4.6), of a real set of data.

Assuming the simulated data, x , were retained from the previous example, the R code to reproduce Figure 4.2 is
$\mathrm{P}=\mathrm{abs}(2 * \mathrm{fft}(\mathrm{x}) / 100)^{\wedge} 2 ; \quad \mathrm{Fr}=0: 99 / 100$ plot(Fr, P, type="o", xlab="frequency", ylab="periodogram")
Different packages scale the FFT differently, so it is a good idea to consult the documentation. R computes it without the factor $n^{-1 / 2}$ and with an additional factor of $e^{2 \pi i \omega_{j}}$ that can be ignored because we will be interested in the squared modulus.

If we consider the data $x_{t}$ in Example 4.1 as a color (waveform) made up of primary colors $x_{t 1}, x_{t 2}, x_{t 3}$ at various strengths (amplitudes), then we might consider the periodogram as a prism that decomposes the color $x_{t}$ into its primary colors (spectrum). Hence the term spectral analysis.

The following is an example using actual data.

## Example 4.3 Star Magnitude

The data in Figure 4.3 are the magnitude of a star taken at midnight for 600 consecutive days. The data are taken from the classic text, The Calculus of Observations, a Treatise on Numerical Mathematics, by E.T. Whittaker and G. Robinson, (1923, Blackie \& Son, Ltd.).

The periodogram for frequencies less than .08 is also displayed in the figure; the periodogram for frequencies higher than .08 are essentially zero. Note that the 29 day cycle and the 25 day cycle are the most prominent periodic components of the data.

The R code to reproduce Figure 4.3 is
$\mathrm{n}=$ length(star)
$\operatorname{par}(\operatorname{mfrow}=c(2,1), \operatorname{mar}=c(3,3,1,1), \operatorname{mgp}=c(1.6, .6,0))$
plot(star, ylab="star magnitude", xlab="day")
Per $=\operatorname{Mod}(f f t(s t a r-m e a n(s t a r)))^{\wedge} 2 / n$
Freq $=(1: n-1) / n$
plot(Freq[1:50], Per[1:50], type='h', lwd=3, ylab="Periodogram", xlab="Frequency")
$\mathrm{u}=$ which.max $(\operatorname{Per}[1: 50]) \quad \# 22$ freq=21/600=.035 cycles/day
$\mathrm{uu}=$ which.max $(\operatorname{Per}[1: 50][-\mathrm{u}])$ \# 25 freq=24/600=.048 cycles/day
1/Freq[22]; 1/Freq[25] \# period = days/cycle
text (.05, 7000, " 25 day cycle"); text(.027, 9000, "29 day cycle")
The periodogram, which was introduced in Schuster (1898) and used in Schuster (1906) for studying the periodicities in the sunspot series (shown in Figure 4.22 in the Problems section) is a sample based statistic. In Example 4.2 and Example 4.3, we discussed the fact that the periodogram may be giving us an idea of the variance components associated with each frequency, as presented in (4.6), of a time series. These variance components, however, are population parameters. The concepts of population parameters and sample statistics, as they relate to spectral analysis of time


Fig. 4.3. Star magnitudes and part of the corresponding periodogram.
series can be generalized to cover stationary time series and that is the topic of the next section.

### 4.3 The Spectral Density

The idea that a time series is composed of periodic components, appearing in proportion to their underlying variances, is fundamental in the spectral representation. The result, called the Spectral Representation Theorem, is quite technical, but the essence of the theorem is that (4.4) is approximately true for any stationary time series.

The examples in the previous section are not generally realistic because time series are rarely the superposition of a small number of distinct sinusoids. An important situation we use repeatedly is the case when the autocovariance function is absolutely summable, in which case there is a spectral density.

## Property 4.1 The Spectral Density

If the autocovariance function, $\gamma(h)$, of a stationary process satisfies

$$
\begin{equation*}
\sum_{h=-\infty}^{\infty}|\gamma(h)|<\infty \tag{4.12}
\end{equation*}
$$

then it has the representation

$$
\begin{equation*}
\gamma(h)=\int_{-1 / 2}^{1 / 2} \mathrm{e}^{2 \pi i \omega h} f(\omega) d \omega \quad h=0, \pm 1, \pm 2, \ldots \tag{4.13}
\end{equation*}
$$

as the inverse transform of the spectral density, which has the representation

$$
\begin{equation*}
f(\omega)=\sum_{h=-\infty}^{\infty} \gamma(h) \mathrm{e}^{-2 \pi i \omega h} \quad-1 / 2 \leq \omega \leq 1 / 2 \tag{4.14}
\end{equation*}
$$

The examples of the previous section were analogues of probability mass functions, or discrete distributions. The pictures of the periodgram in Figure 4.2 and Figure 4.3 are akin to histograms. The spectral density is the analogue of the probability density function, or of continuous distributions.

The fact that $\gamma(h)$ is non-negative definite ensures $f(\omega) \geq 0$ for all $\omega$. It follows immediately from (4.14) that

$$
f(\omega)=f(-\omega)
$$

verifying the spectral density is an even function. Because of the evenness, we will typically only plot $f(\omega)$ for $\omega \geq 0$. In addition, putting $h=0$ in (4.13) yields

$$
\gamma(0)=\operatorname{var}\left(x_{t}\right)=\int_{-1 / 2}^{1 / 2} f(\omega) d \omega
$$

which expresses the total variance as the integrated spectral density over all of the frequencies. We show later on, that a linear filter can isolate the variance in certain frequency intervals or bands.

We note that the absolute summability condition, (4.12), is not satisfied by (4.5), the example that we have used to introduce the idea of a spectral representation. The condition, however, is satisfied for ARMA models.

It is illuminating to examine the spectral density for the series that we have looked at in earlier discussions.

## Example 4.4 White Noise Series

As a simple example, consider the theoretical power spectrum of a sequence of uncorrelated random variables, $w_{t}$, with variance $\sigma_{w}^{2}$. A simulated set of data is displayed in the top of Figure 1.7. Because the autocovariance function was computed in Example 1.14 as $\gamma_{w}(h)=\sigma_{w}^{2}$ for $h=0$, and zero, otherwise, it follows from (4.14), that

$$
f_{w}(\omega)=\sigma_{w}^{2}
$$

for $-1 / 2 \leq \omega \leq 1 / 2$. Hence the process contains equal power at all frequencies. This property is seen in the realization, which seems to contain all different frequencies in a roughly equal mix. In fact, the name white noise comes from the analogy to white light, which contains all frequencies in the color spectrum at the same level of intensity. Figure 4.4 shows a plot of the white noise spectrum for $\sigma_{w}^{2}=1$.

If $x_{t}$ is ARMA, its spectral density can be obtained explicitly using the fact that it is a linear process, i.e., $x_{t}=\sum_{j=0}^{\infty} \psi_{j} w_{t-j}$, where $\sum_{j=0}^{\infty}\left|\psi_{j}\right|<\infty$. In the following property, we exhibit the form of the spectral density of an ARMA model. The proof of the property follows directly from the proof of a more general result, Property 4.5, by using the additional fact that $\psi(z)=\theta(z) / \phi(z)$. The result is analogous to the fact that if $X=a Y$, then $\operatorname{var}(X)=a^{2} \operatorname{var}(Y)$.

## Property 4.2 The Spectral Density of ARMA

If $x_{t}$ is $\operatorname{ARMA}(p, q), \phi(B) x_{t}=\theta(B) w_{t}$, its spectral density is given by

$$
\begin{equation*}
f_{x}(\omega)=\sigma_{w}^{2}\left|\psi\left(\mathrm{e}^{-2 \pi i \omega}\right)\right|^{2}=\sigma_{w}^{2} \frac{\left|\theta\left(\mathrm{e}^{-2 \pi i \omega}\right)\right|^{2}}{\left|\phi\left(\mathrm{e}^{-2 \pi i \omega}\right)\right|^{2}} \tag{4.15}
\end{equation*}
$$

where $\phi(z)=1-\sum_{k=1}^{p} \phi_{k} z^{k}, \theta(z)=1+\sum_{k=1}^{q} \theta_{k} z^{k}$, and $\psi(z)=\sum_{k=0}^{\infty} \psi_{k} z^{k}$.

## Example 4.5 Moving Average

As an example of a series that does not have an equal mix of frequencies, we consider a moving average model. Specifically, consider the MA(1) model given by

$$
x_{t}=w_{t}+.5 w_{t-1} .
$$

A sample realization is shown in the top of Figure 3.2 and we note that the series has less of the higher or faster frequencies. The spectral density will verify this observation.

The autocovariance function is displayed in Example 3.3 on page 58, and for this particular example, we have

$$
\gamma(0)=\left(1+.5^{2}\right) \sigma_{w}^{2}=1.25 \sigma_{w}^{2} ; \quad \gamma( \pm 1)=.5 \sigma_{w}^{2} ; \quad \gamma( \pm h)=0 \text { for } h>1 .
$$

Substituting this directly into the definition given in (4.14), we have

$$
\begin{align*}
f(\omega) & =\sum_{h=-\infty}^{\infty} \gamma(h) \mathrm{e}^{-2 \pi i \omega h}=\sigma_{w}^{2}\left[1.25+.5\left(\mathrm{e}^{-2 \pi i \omega}+\mathrm{e}^{2 \pi i \omega}\right)\right]  \tag{4.16}\\
& =\sigma_{w}^{2}[1.25+\cos (2 \pi \omega)] .
\end{align*}
$$

We can also compute the spectral density using Property 4.2 , which states that for an MA, $f(\omega)=\sigma_{w}^{2}\left|\theta\left(\mathrm{e}^{-2 \pi i \omega}\right)\right|^{2}$. Because $\theta(z)=1+.5 z$, we have

$$
\begin{aligned}
\left|\theta\left(\mathrm{e}^{-2 \pi i \omega}\right)\right|^{2} & =\left|1+.5 \mathrm{e}^{-2 \pi i \omega}\right|^{2}=\left(1+.5 \mathrm{e}^{-2 \pi i \omega}\right)\left(1+.5 \mathrm{e}^{2 \pi i \omega}\right) \\
& =1.25+.5\left(\mathrm{e}^{-2 \pi i \omega}+\mathrm{e}^{2 \pi i \omega}\right)
\end{aligned}
$$

which leads to agreement with (4.16).
Plotting the spectrum for $\sigma_{w}^{2}=1$, as in the middle of Figure 4.4, shows the lower or slower frequencies have greater power than the higher or faster frequencies.

## Example 4.6 A Second-Order Autoregressive Series

We now consider the spectrum of an $\operatorname{AR}(2)$ series of the form

$$
x_{t}-\phi_{1} x_{t-1}-\phi_{2} x_{t-2}=w_{t}
$$

for the special case $\phi_{1}=1$ and $\phi_{2}=-.9$. Figure 1.8 shows a sample realization of such a process for $\sigma_{w}=1$. We note the data exhibit a strong periodic component that makes a cycle about every six points.


Fig. 4.4. Theoretical spectra of white noise (top), a first-order moving average (middle), and a second-order autoregressive process (bottom).

To use Property 4.2, note that $\theta(z)=1, \phi(z)=1-z+.9 z^{2}$ and

$$
\begin{aligned}
\left|\phi\left(e^{-2 \pi i \omega}\right)\right|^{2} & =\left(1-\mathrm{e}^{-2 \pi i \omega}+.9 \mathrm{e}^{-4 \pi i \omega}\right)\left(1-\mathrm{e}^{2 \pi i \omega}+.9 \mathrm{e}^{4 \pi i \omega}\right) \\
& =2.81-1.9\left(\mathrm{e}^{2 \pi i \omega}+\mathrm{e}^{-2 \pi i \omega}\right)+.9\left(\mathrm{e}^{4 \pi i \omega}+\mathrm{e}^{-4 \pi i \omega}\right) \\
& =2.81-3.8 \cos (2 \pi \omega)+1.8 \cos (4 \pi \omega) .
\end{aligned}
$$

Using this result in (4.15), we have that the spectral density of $x_{t}$ is

$$
f_{x}(\omega)=\frac{\sigma_{w}^{2}}{2.81-3.8 \cos (2 \pi \omega)+1.8 \cos (4 \pi \omega)}
$$

Setting $\sigma_{w}=1$, the bottom of Figure 4.4 displays $f_{x}(\omega)$ and shows a strong power component at about $\omega=.16$ cycles per point or a period between six and seven cycles per point and very little power at other frequencies. In this case, modifying the white noise series by applying the second-order AR operator has concentrated the power or variance of the resulting series in a very narrow frequency band.

To reproduce Figure 4.4, use the arma. spec script from astsa (see §R.1): $\operatorname{par}(m f r o w=c(3,1))$
arma.spec (log="no", main="White Noise")
arma.spec (ma=.5, log="no", main="Moving Average")
arma.spec (ar=c (1,-.9), log="no", main="Autoregression")
The above examples motivate the use of the power spectrum for describing the theoretical variance fluctuations of a stationary time series. Indeed, the interpretation
of the spectral density function as the variance of the time series over a given frequency band gives us the intuitive explanation for its physical meaning. The plot of the function $f(\omega)$ over the frequency argument $\omega$ can even be thought of as an analysis of variance, in which the columns or block effects are the frequencies, indexed by $\omega$.

### 4.4 Periodogram and Discrete Fourier Transform

We are now ready to tie together the periodogram, which is the sample-based concept presented in $\S 4.2$, with the spectral density, which is the population-based concept of §4.3.

Definition 4.1 Given data $x_{1}, \ldots, x_{n}$, we define the discrete Fourier transform (DFT) to be

$$
\begin{equation*}
d\left(\omega_{j}\right)=n^{-1 / 2} \sum_{t=1}^{n} x_{t} e^{-2 \pi i \omega_{j} t} \tag{4.17}
\end{equation*}
$$

for $j=0,1, \ldots, n-1$, where the frequencies $\omega_{j}=j / n$ are called the Fourier or fundamental frequencies.

If $n$ is a highly composite integer (i.e., it has many factors), the DFT can be computed by the fast Fourier transform (FFT) introduced in Cooley and Tukey (1965). Sometimes it is helpful to exploit the inversion result for DFTs which shows the linear transformation is one-to-one. For the inverse DFT we have,

$$
\begin{equation*}
x_{t}=n^{-1 / 2} \sum_{j=0}^{n-1} d\left(\omega_{j}\right) \mathrm{e}^{2 \pi i \omega_{j} t} \tag{4.18}
\end{equation*}
$$

for $t=1, \ldots, n$. The following example shows how to calculate the DFT and its inverse in R for the data set $\{1,2,3,4\}$; note that R writes a complex number $z=a+i b$ as a+bi.
(dft = fft(1:4)/sqrt(4))
[1] $5+0 \mathrm{i} \quad-1+1 \mathrm{i}-1+0 \mathrm{i}-1-1 \mathrm{i}$
(idft $=\mathrm{fft}(\mathrm{dft}$, inverse=TRUE)/sqrt(4))
[1] $1+0 \mathrm{i} \quad 2+0 \mathrm{i} \quad 3+0 \mathrm{i} \quad 4+0 \mathrm{i}$
(Re(idft)) \# keep it real
[1] $1 \begin{array}{llll}1 & 3 & 4\end{array}$
We now define the periodogram as the squared modulus ${ }^{3}$ of the DFT.
Definition 4.2 Given data $x_{1}, \ldots, x_{n}$, we define the periodogram to be

$$
\begin{equation*}
I\left(\omega_{j}\right)=\left|d\left(\omega_{j}\right)\right|^{2} \tag{4.19}
\end{equation*}
$$

for $j=0,1,2, \ldots, n-1$.

[^21]Note that $I(0)=n \bar{x}^{2}$, where $\bar{x}$ is the sample mean. Also, for $j \neq 0,4$

$$
\begin{equation*}
I\left(\omega_{j}\right)=\sum_{h=-(n-1)}^{n-1} \widehat{\gamma}(h) e^{-2 \pi i \omega_{j} h} \tag{4.20}
\end{equation*}
$$

In view of (4.20), the periodogram, $I\left(\omega_{j}\right)$, is the sample version of $f\left(\omega_{j}\right)$ given in (4.14). That is, we may think of the periodogram as the "sample spectral density" of $x_{t}$. Although (4.20) seems like a reasonable estimate of $f(\omega)$, recall from Example 4.2 that $I\left(\omega_{j}\right)$, for any $j$, is based on only 2 pieces of information (degrees of freedom).

It is sometimes useful to work with the real and imaginary parts of the DFT individually. To this end, we define the following transforms.

Definition 4.3 Given data $x_{1}, \ldots, x_{n}$, we define the cosine transform

$$
\begin{equation*}
d_{c}\left(\omega_{j}\right)=n^{-1 / 2} \sum_{t=1}^{n} x_{t} \cos \left(2 \pi \omega_{j} t\right) \tag{4.21}
\end{equation*}
$$

and the sine transform

$$
\begin{equation*}
d_{s}\left(\omega_{j}\right)=n^{-1 / 2} \sum_{t=1}^{n} x_{t} \sin \left(2 \pi \omega_{j} t\right) \tag{4.22}
\end{equation*}
$$

where $\omega_{j}=j / n$ for $j=0,1, \ldots, n-1$.
Note that $d_{c}\left(\omega_{j}\right)$ and $d_{s}\left(\omega_{j}\right)$ are averages like the sample mean, but with difference weights (the sample mean has weights $\frac{1}{n}$ for each observation). Under appropriate conditions, there is central limit theorem for these quantities. In non-technical terms, the result is similar to the central limit theorem for sample means, that is,

$$
\begin{equation*}
d_{c}\left(\omega_{j}\right) \dot{\sim} \mathrm{N}\left(0, \frac{1}{2} f\left(\omega_{j}\right)\right) \quad \text { and } \quad d_{s}\left(\omega_{j}\right) \dot{\sim} \mathrm{N}\left(0, \frac{1}{2} f\left(\omega_{j}\right)\right) \tag{4.23}
\end{equation*}
$$

where $\dot{\sim}$ means approximately distributed as for $n$ large. Moreover, it can be shown that for large $n, d_{c}\left(\omega_{j}\right) \perp d_{s}\left(\omega_{j}\right) \perp d_{c}\left(\omega_{k}\right) \perp d_{s}\left(\omega_{k}\right)$, as long as $\omega_{j} \neq \omega_{k}$, where $\perp$ is read is independent of.

We note that $d\left(\omega_{j}\right)=d_{c}\left(\omega_{j}\right)-i d_{s}\left(\omega_{j}\right)$ and hence the periodogram is

$$
\begin{equation*}
I\left(\omega_{j}\right)=d_{c}^{2}\left(\omega_{j}\right)+d_{s}^{2}\left(\omega_{j}\right) \tag{4.24}
\end{equation*}
$$

which for large $n$ is the sum of the squares of two independent normal random variables, which we know has a chi-squared ( $\chi^{2}$ ) distribution. Thus, for large samples, $I\left(\omega_{j}\right) \dot{\sim} \frac{1}{2} f\left(\omega_{j}\right) \chi_{2}^{2}$, or equivalently,

$$
\begin{equation*}
\frac{2 I\left(\omega_{j}\right)}{f\left(\omega_{j}\right)} \dot{\sim} \chi_{2}^{2} \tag{4.27a}
\end{equation*}
$$

[^22]where $\chi_{2}^{2}$ is the chi-squared distribution with 2 degrees of freedom. Since the mean and variance of a $\chi_{v}^{2}$ are $v$ and $2 v$, respectively, it follows from (4.27a) that
\[

$$
\begin{equation*}
E\left[I\left(\omega_{j}\right)\right] \approx f\left(\omega_{j}\right) \quad \text { and } \quad \operatorname{var}\left[I\left(\omega_{j}\right)\right] \approx f^{2}\left(\omega_{j}\right) \tag{4.25}
\end{equation*}
$$

\]

This is bad news because, while the periodgram is approximately unbiased, its variance does not go to zero, and hence it is not consistent. In fact, no matter how large $n$, the variance of the periodogram does not change. Contrast this with the mean $\bar{x}$ of a random sample of size $n$ for which $E[\bar{x}]=\mu$ and $\operatorname{var}[\bar{x}]=\sigma^{2} / n \rightarrow 0$ as $n \rightarrow \infty$.

The technical result regarding the large sample distribution of the periodogram under general conditions is given in the following result.

## Property 4.3 Distribution of the Periodogram Ordinates

If

$$
\begin{equation*}
x_{t}=\sum_{j=-\infty}^{\infty} \psi_{j} w_{t-j}, \quad \sum_{j=-\infty}^{\infty} \sqrt{|j|}\left|\psi_{j}\right|<\infty \tag{4.26}
\end{equation*}
$$

where $w_{t} \sim \operatorname{iid}\left(0, \sigma_{w}^{2}\right)$, then for any collection of $K$ distinct frequencies $\omega_{k} \in(0,1 / 2)$ with $\omega_{k: n} \rightarrow \omega_{k}$ (where $\omega_{k: n}$ is a fundamental frequency) as $n \rightarrow \infty$,

$$
\begin{equation*}
\frac{2 I\left(\omega_{k: n}\right)}{f\left(\omega_{k}\right)} \xrightarrow{d} \operatorname{iid} \chi_{2}^{2} \tag{4.27}
\end{equation*}
$$

provided $f\left(\omega_{k}\right)>0$, for $k=1, \ldots, K$.
The distributional result (4.27) can be used to derive an approximate confidence interval for the spectrum in the usual way. Let $\chi_{v}^{2}(\alpha)$ denote the lower $\alpha$ probability tail for the chi-squared distribution with $v$ degrees of freedom. Then, an approximate $100(1-\alpha) \%$ confidence interval for the spectral density function would be of the form

$$
\begin{equation*}
\frac{2 I\left(\omega_{j: n}\right)}{\chi_{2}^{2}(1-\alpha / 2)} \leq f(\omega) \leq \frac{2 I\left(\omega_{j: n}\right)}{\chi_{2}^{2}(\alpha / 2)} \tag{4.28}
\end{equation*}
$$

The log transform is the variance stabilizing transformation. In this case, the confidence intervals are of the form

$$
\begin{equation*}
\left[\log I\left(\omega_{j: n}\right)+\log 2-\log \chi_{2}^{2}(1-\alpha / 2), \log I\left(\omega_{j: n}\right)+\log 2-\log \chi_{2}^{2}(\alpha / 2)\right] \tag{4.29}
\end{equation*}
$$

Often, nonstationary trends are present that should be eliminated before computing the periodogram. Trends introduce extremely low frequency components in the periodogram that tend to obscure the appearance at higher frequencies. For this reason, it is usually conventional to center the data prior to a spectral analysis using either mean-adjusted data of the form $x_{t}-\bar{x}$ to eliminate the zero or d-c component or to use detrended data of the form $x_{t}-\widehat{\beta}_{1}-\widehat{\beta}_{2} t$. Note that higher order polynomial regressions in $t$ or nonparametric smoothing (linear filtering) could be used in cases where the trend is nonlinear.

As previously indicated, it is often convenient to calculate the DFTs, and hence the periodogram, using the fast Fourier transform algorithm. The FFT utilizes a number


Fig. 4.5. Periodogram of SOI and Recruitment, $n=453\left(n^{\prime}=480\right)$, where the frequency axis is labeled in multiples of $\Delta=1 / 12$. Note the common peaks at $\omega=1 \Delta=1 / 12$, or one cycle per year ( 12 months), and $\omega=\frac{1}{4} \Delta=1 / 48$, or one cycle every four years ( 48 months).
of redundancies in the calculation of the DFT when $n$ is highly composite; that is, an integer with many factors of 2,3 , or 5 , the best case being when $n=2^{p}$ is a factor of 2. Details may be found in Cooley and Tukey (1965). To accommodate this property, we can pad the centered (or detrended) data of length $n$ to the next highly composite integer $n^{\prime}$ by adding zeros, i.e., setting $x_{n+1}^{c}=x_{n+2}^{c}=\cdots=x_{n^{\prime}}^{c}=0$, where $x_{t}^{c}$ denotes the centered data. This means that the fundamental frequency ordinates will be $\omega_{j}=j / n^{\prime}$ instead of $j / n$. We illustrate by considering the periodogram of the SOI and Recruitment series, as has been given in Figure 1.5 of Chapter 1. Recall that they are monthly series and $n=453$ months. To find $n^{\prime}$ in R , use the command nextn (453) to see that $n^{\prime}=480$ will be used in the spectral analyses by default.

## Example 4.7 Periodogram of SOI and Recruitment Series

Figure 4.5 shows the periodograms of each series, where the frequency axis is labeled in multiples of $\Delta=1 / 12$. As previously indicated, the centered data have been padded to a series of length 480 . We notice a narrow-band peak at the obvious yearly ( 12 month) cycle, $\omega=1 \Delta=1 / 12$. In addition, there is considerable power in a wide band at the lower frequencies that is centered around the four-year (48 month) cycle $\omega=\frac{1}{4} \Delta=1 / 48$ representing a possible El Niño effect. This wide


Fig. 4.6. Log-periodogram of SOI and Recruitment. 95\% confidence intervals are indicated by the blue line in the upper right corner.
band activity suggests that the possible El Niño cycle is irregular, but tends to be around four years on average. We will continue to address this problem as we move to more sophisticated analyses.

Noting $\chi_{2}^{2}(.025)=.05$ and $\chi_{2}^{2}(.975)=7.38$, we can obtain approximate $95 \%$ confidence intervals for the frequencies of interest. For example, the periodogram of the SOI series is $I_{S}(1 / 12)=.97$ at the yearly cycle. An approximate $95 \%$ confidence interval for the spectrum $f_{S}(1 / 12)$ is then

$$
[2(.97) / 7.38,2(.97) / .05]=[.26,38.4]
$$

which is too wide to be of much use. We do notice, however, that the lower value of .26 is higher than any other periodogram ordinate, so it is safe to say that this value is significant. On the other hand, an approximate $95 \%$ confidence interval for the spectrum at the four-year cycle, $f_{S}(1 / 48)$, is

$$
[2(.05) / 7.38,2(.05) / .05]=[.01,2.12]
$$

which again is extremely wide, and with which we are unable to establish significance of the peak.

We now give the R commands that can be used to reproduce Figure 4.5. To calculate and graph the periodogram, we used the mvspec script from astsa, although

R's spec.pgram can be used. In the code, we set log="no" because the script will plot the periodogram on a $\log _{10}$ scale by default. Figure 4.5 displays a bandwidth, which we discuss in the next section.

```
require(astsa) # needed for mvspec() - otherwise use spec.pgram()
par(mfrow=c(2,1))
soi.per = mvspec(soi, log="no"); abline(v=1/4, lty="dotted")
rec.per = mvspec(rec, log="no"); abline(v=1/4, lty="dotted")
```

The confidence intervals for the SOI series at the yearly cycle, $\omega=1 / 12=$ $40 / 480$, and the possible El Niño cycle of four years $\omega=1 / 48=10 / 480$ can be computed in R as follows:

```
soi.per$spec[40] # 0.97223
soi pgram at freq 1/12 = 40/480
soi.per$spec[10] # 0.05372
soi pgram at freq 1/48 = 10/480
# conf intervals - returned value:
U = qchisq(.025,2) # 0.05063
L = qchisq(.975,2) # 7.37775
2*soi.per$spec[10]/L # 0.01456
2*soi.per$spec[10]/U # 2.12220
2*soi.per$spec[40]/L # 0.26355
2*soi.per$spec[40]/U # 38.40108
```

The example above makes it clear that the periodogram as an estimator is susceptible to large uncertainties, and we need to find a way to reduce the variance. Not surprisingly, this result follows if we think about the periodogram, $I\left(\omega_{j}\right)$ as an estimator of the spectral density $f(\omega)$ and realize that it is the sum of squares of only two random variables for any sample size. The solution to this dilemma is suggested by the analogy with classical statistics where we look for independent random variables with the same variance and average the squares of these common variance observations.

As an analogy to using the periodogram to estimate the spectral density, consider the problem of taking a random sample and then trying to estimate a probability density based on a histogram with many cells. This approach is demonstrated in Figure 4.7.

### 4.5 Nonparametric Spectral Estimation

To continue the discussion that ended the previous section, we introduce a frequency band, $\mathcal{B}$, of $L \ll n$ contiguous fundamental frequencies, centered around frequency $\omega_{j}=j / n$, which is chosen close to a frequency of interest, $\omega$. Let

$$
\begin{equation*}
\mathcal{B}=\left\{\omega_{j}+k / n: k=0, \pm 1, \ldots, \pm m\right\}, \tag{4.30}
\end{equation*}
$$

where

$$
\begin{equation*}
L=2 m+1 \tag{4.31}
\end{equation*}
$$

is an odd number, chosen such that the spectral values in the interval $\mathcal{B}$,

$$
f\left(\omega_{j}+k / n\right), \quad k=-m, \ldots, 0, \ldots, m
$$



Fig. 4.7. Histogram of a sample of $n=200$ standard normals with 100 cells and with the standard normal density superimposed. The periodogram is to the spectral density as the histogram is to the normal density.
are approximately equal to $f(\omega)$. For example, to see a small section of the $\operatorname{AR}(2)$ spectrum|near the peak|shown in Figure 4.4, use
arma.spec (ar=c (1,-.9), xlim=c (.15,.151))
which is displayed in Figure 4.8.
We now define an averaged (or smoothed) periodogram as the average of the periodogram values, say,

$$
\begin{equation*}
\bar{f}(\omega)=\frac{1}{L} \sum_{k=-m}^{m} I\left(\omega_{j}+k / n\right) \tag{4.32}
\end{equation*}
$$

over the band $\mathcal{B}$. Under the assumption that the spectral density is fairly constant in the band $\mathcal{B}$, and in view of (4.27) we can show that under appropriate conditions, for large $n$, the periodograms in (4.32) are approximately distributed as independent $f(\omega) \chi_{2}^{2} / 2$ random variables, for $0<\omega<1 / 2$, as long as we keep $L$ fairly small relative to $n$. Thus, under these conditions, $L \bar{f}(\omega)$ is the sum of $L$ approximately independent $f(\omega) \chi_{2}^{2} / 2$ random variables. It follows that, for large $n$,

$$
\begin{equation*}
\frac{2 L \bar{f}(\omega)}{f(\omega)} \dot{\sim} \chi_{2 L}^{2} \tag{4.33}
\end{equation*}
$$

Now we have

$$
\begin{equation*}
E[\bar{f}(\omega)] \approx f(\omega) \quad \text { and } \quad \operatorname{var}[\bar{f}(\omega)] \approx f^{2}(\omega) / L \tag{4.34}
\end{equation*}
$$

which can be compared to (4.25). In this case, we have consistency if we let $L \rightarrow \infty$ as $n \rightarrow \infty$, but $L$ must grow much slower than $n$, of course (in fact, $L / n \rightarrow 0$ as $n \rightarrow \infty)$.

In this scenario, where we smooth the periodogram by simple averaging, the width of the frequency interval defined by (4.30),

$$
\begin{equation*}
B_{w}=\frac{L}{n} \tag{4.35}
\end{equation*}
$$

is called the bandwidth. Note (4.35) implies the degrees of freedom can be expressed as


Fig. 4.8. A small section (near the peak) of the $\operatorname{AR}(2)$ spectrum shown in Figure 4.4 .

$$
\begin{equation*}
2 L=2 B_{w} n, \tag{4.36}
\end{equation*}
$$

or twice the time-bandwidth product. The result (4.33) can be rearranged to obtain an approximate $100(1-\alpha) \%$ confidence interval of the form

$$
\begin{equation*}
\frac{2 L \bar{f}(\omega)}{\chi_{2 L}^{2}(1-\alpha / 2)} \leq f(\omega) \leq \frac{2 L \bar{f}(\omega)}{\chi_{2 L}^{2}(\alpha / 2)} \tag{4.37}
\end{equation*}
$$

for the true spectrum, $f(\omega)$.
As previously discussed, the visual impact of a spectral density plot will be improved by plotting the logarithm of the spectrum, which is the variance stabilizing transformation in this situation. This phenomenon can occur when regions of the spectrum exist with peaks of interest much smaller than some of the main power components. For the log spectrum, we obtain an interval of the form

$$
\begin{align*}
& {\left[\log \bar{f}(\omega)+\log 2 L-\log \chi_{2 L}^{2}(1-\alpha / 2)\right.} \\
& \left.\log \bar{f}(\omega)+\log 2 L-\log \chi_{2 L}^{2}(\alpha / 2)\right] \tag{4.38}
\end{align*}
$$

If zeros are appended before computing the spectral estimators, we need to adjust the degrees of freedom and an approximation is to replace $2 L$ by $2 L n / n^{\prime}$. Hence, we define the adjusted degrees of freedom as

$$
\begin{equation*}
d f=\frac{2 L n}{n^{\prime}} \tag{4.39}
\end{equation*}
$$

and use it instead of $2 L$ in the confidence intervals (4.37) and (4.38). For example, (4.37) becomes

$$
\begin{equation*}
\frac{d f \bar{f}(\omega)}{\chi_{d f}^{2}(1-\alpha / 2)} \leq f(\omega) \leq \frac{d f \bar{f}(\omega)}{\chi_{d f}^{2}(\alpha / 2)} \tag{4.40}
\end{equation*}
$$

Before proceeding further, we pause to consider computing the average periodograms for the SOI and Recruitment series, as shown in Figure 4.9.

## Example 4.8 Averaged Periodogram for SOI and Recruitment

Generally, it is a good idea to try several bandwidths that seem to be compatible with the general overall shape of the spectrum, as suggested by the periodogram. The SOI and Recruitment series periodograms, previously computed in Figure 4.5, suggest the power in the lower El Niño frequency needs smoothing to identify the


Fig. 4.9. The averaged periodogram of the SOI and Recruitment series $n=453, n^{\prime}=480, L=$ $9, d f=17$, showing common peaks at the four year period, $\omega=\frac{1}{4} \Delta=1 / 48$ (cycle/months), the yearly period, $\omega=1 \Delta=1 / 12$ (cycle/months) and some of its harmonics $\omega=k \Delta$ for $k=2,3$.


Fig. 4.10. Figure 4.9 with the average periodogram ordinates plotted on a $\log _{10}$ scale. The display in the upper right-hand corner represents a generic $95 \%$ confidence interval.
predominant overall period. Trying values of $L$ leads to the choice $L=9$ as a reasonable value, and the result is displayed in Figure 4.9.

The smoothed spectra shown in Figure 4.9 provide a sensible compromise between the noisy version, shown in Figure 4.5, and a more heavily smoothed spectrum, which might lose some of the peaks. An undesirable effect of averaging can be noticed at the yearly cycle, $\omega=1 \Delta$, where the narrow band peaks that appeared in the periodograms in Figure 4.5 have been flattened and spread out to nearby frequencies. We also notice, and have marked, the appearance of harmonics of the yearly cycle, that is, frequencies of the form $\omega=k \Delta$ for $k=1,2, \ldots$. Harmonics typically occur when a periodic component is present, but not in a sinusoidal fashion; see Example 4.9.

Figure 4.9 can be reproduced in R using the following commands. The basic call is to the function mvspec, which is available in astsa; alternately, use R's spec. pgram. To compute averaged periodograms, use the Daniell kernel, and specify $m$, where $L=2 m+1$ ( $L=9$ and $m=4$ in this example). We will explain the kernel concept later in this section, specifically just prior to Example 4.10.

```
par(mfrow=c (2,1))
(k = kernel("daniell", 4))
soi.ave = mvspec(soi, k, log="no")
abline(v=c(.25,1,2,3), lty=2)
# Repeat above lines using rec in place of soi on line 3
soi.ave$bandwidth # = 0.225
```

The displayed bandwidth (.225) is adjusted for the fact that the frequency scale of the plot is in terms of cycles per year instead of cycles per month (the original unit of the data). Using (4.35), the bandwidth in terms of months is $9 / 480=.01875$; the displayed value is simply converted to years, $.01875 \frac{\text { cycles }}{\text { month }} \times 12 \frac{\text { months }}{\text { year }}=.225 \frac{\text { cycles }}{\text { year }}$.

The adjusted degrees of freedom are $d f=2(9)(453) / 480 \approx 17$. We can use this value for the $95 \%$ confidence intervals, with $\chi_{d f}^{2}(.025)=7.56$ and $\chi_{d f}^{2}(.975)=$ 30.17. Substituting into (4.40) gives the intervals in Table 4.1 for the two frequency bands identified as having the maximum power. To examine the two peak power possibilities, we may look at the $95 \%$ confidence intervals and see whether the lower limits are substantially larger than adjacent baseline spectral levels. For example, the El Niño frequency of 48 months has lower limits that exceed the values the spectrum would have if there were simply a smooth underlying spectral function without the peaks. The relative distribution of power over frequencies is different, with the SOI having less power at the lower frequency, relative to the seasonal periods, and the recruit series having relatively more power at the lower or El Niño frequency.

The entries in Table 4.1 for SOI can be obtained in R as follows:

```
df = soi.ave$df # df = 16.9875 (returned values)
U = qchisq(.025, df) # U = 7.555916
L = qchisq(.975, df) # L = 30.17425
soi.ave$spec[10] # 0.0495202
soi.ave$spec[40] # 0.1190800
# intervals
df*soi.ave$spec[10]/L # Q.0278789
df*soi.ave$spec[10]/U # 0.1113333
df*soi.ave$spec[40]/L # 0.0670396
df*soi.ave$spec[40]/U # 0.2677201
# repeat above commands with soi replaced by rec
```

Table 4.1. Confidence Intervals for the Spectra of the SOI and Recruitment Series

| Series | $\omega$ | Period | Power | Lower | Upper |
| :---: | :---: | :---: | :---: | :---: | :---: |
| SOI | $1 / 48$ | 4 years | .05 | .03 | .11 |
|  | $1 / 12$ | 1 year | .12 | .07 | .27 |
| Recruits | $1 / 48$ | 4 years | 6.59 | 3.71 | 14.82 |
| $\times 10^{2}$ | $1 / 12$ | 1 year | 2.19 | 1.24 | 4.93 |

Finally, Figure 4.10 shows the averaged periodograms in Figure 4.9 plotted on a $\log _{10}$ scale. This is the default plot in R, and these graphs can be obtained by removing the statement log="no". Notice that the default plot also shows a generic confidence interval of the form (4.38) in the upper right-hand corner. To use it, imagine placing the tick mark on the averaged periodogram ordinate of interest; the resulting bar then constitutes an approximate $95 \%$ confidence interval for the spectrum at that frequency. We note that displaying the estimates on a log scale tends to emphasize the harmonic components.

## Example 4.9 Harmonics

In the previous example, we saw that the spectra of the annual signals displayed minor peaks at the harmonics; that is, the signal spectra had a large peak at $\omega=$ $1 \Delta=1 / 12$ cycles/month (the one-year cycle) and minor peaks at its harmonics $\omega=k \Delta$ for $k=2,3, \ldots$ (two-, three-, and so on, cycles per year). This will often be the case because most signals are not perfect sinusoids (or perfectly cyclic). In this case, the harmonics are needed to capture the non-sinusoidal behavior of the signal. As an example, consider the signal formed in Figure 4.11 from a (fundamental) sinusoid oscillating at two cycles per unit time along with the second through sixth harmonics at decreasing amplitudes. In particular, the signal was formed as

$$
\begin{align*}
x_{t}= & \sin (2 \pi 2 t)+.5 \sin (2 \pi 4 t)+.4 \sin (2 \pi 6 t) \\
& +.3 \sin (2 \pi 8 t)+.2 \sin (2 \pi 10 t)+.1 \sin (2 \pi 12 t) \tag{4.41}
\end{align*}
$$

for $0 \leq t \leq 1$. Notice that the signal is non-sinusoidal in appearance and rises quickly then falls slowly.

A figure similar to Figure 4.11 can be generated in R as follows.

```
t = seq(0, 1, by=1/200)
amps = c(1, .5, .4, .3, .2, .1)
x = matrix(0, 201, 6)
for (j in 1:6) x[,j] = amps[j]*sin(2*pi*t*2*j)
x = ts(cbind(x, rowSums(x)), start=0, deltat=1/200)
ts.plot(x, lty=c(1:6, 1), lwd=c(rep(1,6), 2), ylab="Sinusoids")
names = c("Fundamental","2nd Harmonic","3rd Harmonic","4th Harmonic", "5th
    Harmonic", "6th Harmonic", "Formed Signal")
legend("topright", names, lty=c(1:6, 1), lwd=c(rep(1,6), 2))
```

Example 4.8 points out the necessity for having some relatively systematic procedure for deciding whether peaks are significant. The question of deciding whether a single peak is significant usually rests on establishing what we might think of as


Fig. 4.11. A signal (thick solid line) formed by a fundamental sinusoid (thin solid line) oscillating at two cycles per unit time and its harmonics as specified in (4.41).
a baseline level for the spectrum, defined rather loosely as the shape that one would expect to see if no spectral peaks were present. This profile can usually be guessed by looking at the overall shape of the spectrum that includes the peaks; usually, a kind of baseline level will be apparent, with the peaks seeming to emerge from this baseline level. If the lower confidence limit for the spectral value is still greater than the baseline level at some predetermined level of significance, we may claim that frequency value as a statistically significant peak. To be consistent with our stated indifference to the upper limits, we might use a one-sided confidence interval.

Care must be taken when we make a decision about the bandwidth $B_{w}$ over which the spectrum will be essentially constant. Taking too broad a band will tend to smooth out valid peaks in the data when the constant variance assumption is not met over the band. Taking too narrow a band will lead to confidence intervals so wide that peaks are no longer statistically significant. Thus, we note that there is a conflict here between variance properties or bandwidth stability, which can be improved by increasing $B_{w}$ and resolution, which can be improved by decreasing $B_{w}$. A common approach is to try a number of different bandwidths and to look qualitatively at the spectral estimators for each case.

To address the problem of resolution, it should be evident that the flattening of the peaks in Figure 4.9 and Figure 4.10 was due to the fact that simple averaging was used in computing $\bar{f}(\omega)$ defined in (4.32). There is no particular reason to use simple averaging, and we might improve the estimator by employing a weighted average, say

$$
\begin{equation*}
\widehat{f}(\omega)=\sum_{k=-m}^{m} h_{k} I\left(\omega_{j}+k / n\right) \tag{4.42}
\end{equation*}
$$

using the same definitions as in (4.32) but where the weights $h_{k}>0$ satisfy

$$
\sum_{k=-m}^{m} h_{k}=1
$$

In particular, it seems reasonable that the resolution of the estimator will improve if we use weights that decrease as distance from the center weight $h_{0}$ increases; we will return to this idea shortly. To obtain the averaged periodogram, $\bar{f}(\omega)$, in (4.42), set $h_{k}=L^{-1}$, for all $k$, where $L=2 m+1$. The asymptotic theory established for $\bar{f}(\omega)$ still holds for $\widehat{f}(\omega)$ provided that the weights satisfy the additional condition that if $m \rightarrow \infty$ as $n \rightarrow \infty$ but $m / n \rightarrow 0$, then

$$
\sum_{k=-m}^{m} h_{k}^{2} \rightarrow 0
$$

Under these conditions, for $n$ large, we have

$$
\begin{equation*}
E[\widehat{f}(\omega)] \approx f(\omega) \quad \text { and } \quad \operatorname{var}[\widehat{f}(\omega)] \approx f^{2}(\omega) \sum_{k=-m}^{m} h_{k}^{2} \tag{4.43}
\end{equation*}
$$

which can be compared to (4.34); as before, we have that $\widehat{f}(\omega)$ is consistent. We have already seen this result in the case of $\bar{f}(\omega)$, where the weights are constant, $h_{k}=L^{-1}$, in which case $\sum_{k=-m}^{m} h_{k}^{2}=L^{-1}$. The distributional properties of (4.42) are more difficult now because $\widehat{f}(\omega)$ is a weighted linear combination of asymptotically independent $\chi^{2}$ random variables. An approximation that seems to work well is to replace $L$ by $\left(\sum_{k=-m}^{m} h_{k}^{2}\right)^{-1}$. That is, define

$$
\begin{equation*}
L_{h}=\left(\sum_{k=-m}^{m} h_{k}^{2}\right)^{-1} \tag{4.44}
\end{equation*}
$$

and use the approximation

$$
\begin{equation*}
\frac{2 L_{h} \widehat{f}(\omega)}{f(\omega)} \dot{\sim} \chi_{2 L_{h}}^{2} \tag{4.45}
\end{equation*}
$$

In analogy to (4.35), we will define the bandwidth in this case to be

$$
\begin{equation*}
B_{w}=\frac{L_{h}}{n} \tag{4.46}
\end{equation*}
$$

Using the approximation (4.45) we obtain an approximate $100(1-\alpha) \%$ confidence interval of the form

$$
\begin{equation*}
\frac{2 L_{h} \widehat{f}(\omega)}{\chi_{2 L_{h}}^{2}(1-\alpha / 2)} \leq f(\omega) \leq \frac{2 L_{h} \widehat{f}(\omega)}{\chi_{2 L_{h}}^{2}(\alpha / 2)} \tag{4.47}
\end{equation*}
$$

for the true spectrum, $f(\omega)$. If the data are padded to $n^{\prime}$, then replace $2 L_{h}$ in (4.47) with $d f=2 L_{h} n / n^{\prime}$ as in (4.39).


Fig. 4.12. Modified Daniell kernel weights used in Example 4.10

An easy way to generate the weights in R is by repeated use of the Daniell kernel. For example, with $m=1$ and $L=2 m+1=3$, the Daniell kernel has weights $\left\{h_{k}\right\}=\left\{\frac{1}{3}, \frac{1}{3}, \frac{1}{3}\right\}$; applying this kernel to a sequence of numbers, $\left\{u_{t}\right\}$, produces

$$
\widehat{u}_{t}=\frac{1}{3} u_{t-1}+\frac{1}{3} u_{t}+\frac{1}{3} u_{t+1} .
$$

We can apply the same kernel again to the $\widehat{u}_{t}$,

$$
\widehat{\hat{u}_{t}}=\frac{1}{3} \widehat{u}_{t-1}+\frac{1}{3} \widehat{u}_{t}+\frac{1}{3} \widehat{u}_{t+1}
$$

which simplifies to

$$
\widehat{\widehat{u}}_{t}=\frac{1}{9} u_{t-2}+\frac{2}{9} u_{t-1}+\frac{3}{9} u_{t}+\frac{2}{9} u_{t+1}+\frac{1}{9} u_{t+2} .
$$

The modified Daniell kernel puts half weights at the end points, so with $m=1$ the weights are $\left\{h_{k}\right\}=\left\{\frac{1}{4}, \frac{2}{4}, \frac{1}{4}\right\}$ and

$$
\widehat{u}_{t}=\frac{1}{4} u_{t-1}+\frac{1}{2} u_{t}+\frac{1}{4} u_{t+1} .
$$

Applying the same kernel again to $\widehat{u}_{t}$ yields

$$
\widehat{\hat{u}_{t}}=\frac{1}{16} u_{t-2}+\frac{4}{16} u_{t-1}+\frac{6}{16} u_{t}+\frac{4}{16} u_{t+1}+\frac{1}{16} u_{t+2} .
$$

These coefficients can be obtained in R by issuing the kernel command. For example, kernel("modified.daniell", c(1,1)) would produce the coefficients of the last example.

## Example 4.10 Smoothed Periodogram for SOI and Recruitment

In this example, we estimate the spectra of the SOI and Recruitment series using the smoothed periodogram estimate in (4.42). We used a modified Daniell kernel twice, with $m=3$ both times. This yields $L_{h}=1 / \sum h_{k}^{2}=9.232$, which is close to the value of $L=9$ used in Example 4.8. In this case, the bandwidth is $B_{w}=9.232 / 480=.019$ and the modified degrees of freedom is $d f=2 L_{h} 453 / 480=17.43$. The weights, $h_{k}$, can be obtained and graphed in R as follows; see Figure 4.12.
$\begin{array}{ll}\text { kernel("modified.daniell", c(3,3)) } & \text { \# for a list } \\ \text { plot(kernel("modified.daniell", c(3,3))) } & \text { \# for a plot }\end{array}$ plot(kernel("modified.daniell", c(3,3))) \# for a plot

Series: soi
Smoothed Periodogram


Series: soi
Smoothed Periodogram


Series: rec
Smoothed Periodogram


Fig. 4.13. Smoothed spectral estimates of the SOI and Recruitment series; see Example 4.10 for details. The bottom pair are drawn on a $\log$ scale. Note that the El Niño cycle is a broad band cycle ranging from about 3 to 7 years. See Figure 4.18 for a graph of the approximate El Niño cycle.

The resulting spectral estimates can be viewed in Figure 4.13 and we notice that the estimates more appealing than those in Figure 4.9. Figure 4.13 was generated in R as follows; we also show how to obtain $d f$ and $B_{w}$.

```
par(mfrow=c (2,1))
k = kernel('modified.daniell", c(3,3))
soi.smo = mvspec(soi, k, log="no")
abline(v=1, lty="dotted"); abline(v=1/4, lty="dotted")
## Repeat above lines with rec replacing soi in line 3
df = soi.smo$df # df = 17.42618
soi.smo$bandwidth # Bw = 0.2308103 = 12*9.232/480
```

Reissuing the mvspec commands with log="no" removed will result in a figure similar to Figure 4.10; see Figure 4.13. Finally, we mention that the modified Daniell kernel is used by default. For example, an easier way to obtain soi.smo is to issue the command:

```
soi.smo = mvspec(soi, spans=c(7,7))
```

Notice that spans is a vector of odd integers, given in terms of $L=2 m+1$ instead of $m$. These values give the widths of the modified Daniell smoother to be used to smooth the periodogram.

## TApering

We are now ready to briefly introduce the concept of tapering; a more detailed discussion may be found in Bloomfield (2000, §9.5). Suppose $x_{t}$ is a mean-zero, stationary process with spectral density $f_{x}(\omega)$. If we replace the original series by the tapered series

$$
\begin{equation*}
y_{t}=h_{t} x_{t} \tag{4.48}
\end{equation*}
$$

for $t=1,2, \ldots, n$, use the modified DFT

$$
\begin{equation*}
d_{y}\left(\omega_{j}\right)=n^{-1 / 2} \sum_{t=1}^{n} h_{t} x_{t} \mathrm{e}^{-2 \pi i \omega_{j} t} \tag{4.49}
\end{equation*}
$$

and let $I_{y}\left(\omega_{j}\right)=\left|d_{y}\left(\omega_{j}\right)\right|^{2}$, we will obtain

$$
\begin{equation*}
E\left[I_{y}\left(\omega_{j}\right)\right]=\int_{-1 / 2}^{1 / 2} W_{n}\left(\omega_{j}-\omega\right) f_{x}(\omega) d \omega \tag{4.50}
\end{equation*}
$$

The value $W_{n}(\omega)$ is called a spectral window because, in view of (4.50), it is determining which part of the spectral density $f_{x}(\omega)$ is being "seen" by the estimator $I_{y}\left(\omega_{j}\right)$ on average. In the case that $h_{t}=1$ for all $t, I_{y}\left(\omega_{j}\right)=I_{x}\left(\omega_{j}\right)$ is simply the periodogram of the data and the window is

$$
\begin{equation*}
W_{n}(\omega)=\frac{\sin ^{2}(n \pi \omega)}{n \sin ^{2}(\pi \omega)} \tag{4.51}
\end{equation*}
$$

with $W_{n}(0)=n$, which is known as the Fejér or modified Bartlett kernel. If we consider the averaged periodogram in (4.32), namely

$$
\bar{f}_{x}(\omega)=\frac{1}{L} \sum_{k=-m}^{m} I_{x}\left(\omega_{j}+k / n\right)
$$



Fig. 4.14. Averaged Fejér window (top row) and the corresponding cosine taper window (bottom row) for $L=9, n=480$. The extra tic marks on the horizontal axis of the left-hand plots exhibit the predicted bandwidth, $B_{w}=9 / 480=.01875$.
the window, $W_{n}(\omega)$, in (4.50) will take the form

$$
\begin{equation*}
W_{n}(\omega)=\frac{1}{n L} \sum_{k=-m}^{m} \frac{\sin ^{2}[n \pi(\omega+k / n)]}{\sin ^{2}[\pi(\omega+k / n)]} . \tag{4.52}
\end{equation*}
$$

Tapers generally have a shape that enhances the center of the data relative to the extremities, such as a cosine bell of the form

$$
\begin{equation*}
h_{t}=.5\left[1+\cos \left(\frac{2 \pi(t-\bar{t})}{n}\right)\right], \tag{4.53}
\end{equation*}
$$

where $\bar{t}=(n+1) / 2$, favored by Blackman and Tukey (1959). In Figure 4.14, we have plotted the shapes of two windows, $W_{n}(\omega)$, for $n=480$ and $L=9$, when (i) $h_{t} \equiv 1$, in which case, (4.52) applies, and (ii) $h_{t}$ is the cosine taper in (4.53). In both cases the predicted bandwidth should be $B_{w}=9 / 480=.01875$ cycles per point, which corresponds to the "width" of the windows shown in Figure 4.14. Both windows produce an integrated average spectrum over this band but the untapered window in the top panels shows considerable ripples over the band and outside the band. The ripples outside the band are called sidelobes and tend to introduce frequencies from outside the interval that may contaminate the desired spectral estimate within the


Fig. 4.15. Smoothed spectral estimates of the SOI without tapering (dashed line) and with full tapering (solid line); see Example 4.11 for details.
band. For example, a large dynamic range for the values in the spectrum introduces spectra in contiguous frequency intervals several orders of magnitude greater than the value in the interval of interest. This effect is sometimes called leakage. Figure 4.14 emphasizes the suppression of the sidelobes in the Fejér kernel when a cosine taper is used.

## Example 4.11 The Effect of Tapering the SOI Series

In this example, we examine the effect of tapering on the estimate of the spectrum of the SOI series. The results for the Recruitment series are similar. Figure 4.15 shows two spectral estimates plotted on a log scale. The degree of smoothing here is the same as in Example 4.10. The dashed line in Figure 4.15 shows the estimate without any tapering and hence it is the same as the estimated spectrum displayed in the top of Figure 4.13. The solid line shows the result with full tapering. Notice that the tapered spectrum does a better job in separating the yearly cycle $(\omega=1)$ and the El Niño cycle ( $\omega=1 / 4$ ).

The following R session was used to generate Figure 4.15 . We note that, by default, mvspec does not taper. For full tapering, we use the argument taper=. 5 to instruct mvspec to taper $50 \%$ of each end of the data; any value between 0 and .5 is acceptable.

```
s0 = mvspec(soi, spans=c(7,7), plot=FALSE) # no taper
s50 = mvspec(soi, spans=c(7,7), taper=.5, plot=FALSE) # full taper
plot(s0$freq, s0$spec, log="y", type="l", lty=2, ylab="spectrum",
    xlab="frequency") # dashed line
lines(s50$freq, s50$spec) # solid line
```


### 4.6 Parametric Spectral Estimation

The methods of $\S 4.5$ lead to estimators generally referred to as nonparametric spectra because no assumption is made about the parametric form of the spectral density. In Property 4.2, we exhibited the spectrum of an ARMA process and we might consider basing a spectral estimator on this function, substituting the parameter estimates from an $\operatorname{ARMA}(p, q)$ fit on the data into the formula for the spectral density $f_{x}(\omega)$ given in (4.15). Such an estimator is called a parametric spectral estimator. For convenience, a parametric spectral estimator is obtained by fitting an $\operatorname{AR}(p)$ to the data, where the order $p$ is determined by one of the model selection criteria, such as AIC, AICc, and BIC, defined in (2.15)-(2.17). Parametric autoregressive spectral estimators will often have superior resolution in problems when several closely spaced narrow spectral peaks are present and are preferred by engineers for a broad variety of problems (see Kay, 1988). The development of autoregressive spectral estimators has been summarized by Parzen (1983).

If $\widehat{\phi}_{1}, \widehat{\phi}_{2}, \ldots, \widehat{\phi}_{p}$ and $\widehat{\sigma}_{w}^{2}$ are the estimates from an $\operatorname{AR}(p)$ fit to $x_{t}$, then based on Property 4.2, a parametric spectral estimate of $f_{x}(\omega)$ is attained by substituting these estimates into (4.15), that is,

$$
\begin{equation*}
\widehat{f}_{x}(\omega)=\frac{\widehat{\sigma}_{w}^{2}}{\left|\widehat{\phi}\left(\mathrm{e}^{-2 \pi i \omega}\right)\right|^{2}} \tag{4.54}
\end{equation*}
$$

where

$$
\begin{equation*}
\widehat{\phi}(z)=1-\widehat{\phi}_{1} z-\widehat{\phi}_{2} z^{2}-\cdots-\widehat{\phi}_{p} z^{p} . \tag{4.55}
\end{equation*}
$$

An interesting fact about rational spectra of the form (4.15) is that any spectral density can be approximated, arbitrarily close, by the spectrum of an AR process.

## Property 4.4 AR Spectral Approximation

Let $g(\omega)$ be the spectral density of a stationary process, $x_{t}$. Then, given $\epsilon>0$, there is an $A R(p)$ representation

$$
x_{t}=\sum_{k=1}^{p} \phi_{k} x_{t-k}+w_{t}
$$

where $w_{t}$ is white noise with variance $\sigma_{w}^{2}$, such that

$$
\left|f_{x}(\omega)-g(\omega)\right|<\epsilon \quad \text { for all } \omega \in[-1 / 2,1 / 2]
$$

Moreover, $p$ is finite and the roots of $\phi(z)=1-\sum_{k=1}^{p} \phi_{k} z^{k}$ are outside the unit circle.
One drawback, however, is that the property does not tell us how large $p$ must be before the approximation is reasonable; in some situations $p$ may be extremely large. Property 4.4 also holds for MA and for ARMA processes in general, and a proof of the result may be found in Fuller (1996, Ch 4). We demonstrate the technique in the following example.


Fig. 4.16. Model selection criteria AIC and BIC as a function of order $p$ for autoregressive models fitted to the SOI series.

## Example 4.12 Autoregressive Spectral Estimator for SOI

Consider obtaining results comparable to the nonparametric estimators shown in Figure 4.9 for the SOI series. Fitting successively higher order $\operatorname{AR}(p)$ models for $p=1,2, \ldots, 30$ yields a minimum BIC and a minimum AIC at $p=15$, as shown in Figure 4.16. We can see from Figure 4.16 that BIC is very definite about which model it chooses; that is, the minimum BIC is very distinct. On the other hand, it is not clear what is going to happen with AIC; that is, the minimum is not so clear, and there is some concern that AIC will start decreasing after $p=30$. Minimum AICc selects the $p=15$ model, but suffers from the same uncertainty as AIC. The spectrum is shown in Figure 4.17, and we note the strong peaks at 52 months and 12 months corresponding to the nonparametric estimators obtained in $\S 4.5$. In addition, the harmonics of the yearly period are evident in the estimated spectrum.

To perform a similar analysis in $R$, the command spec. ar can be used to fit the best model via AIC and plot the resulting spectrum. A quick way to obtain the AIC values is to run the ar command as follows.

```
spaic = spec.ar(soi, log="no") # min AIC spec
abline(v=frequency(soi)*1/52, lty="dotted") # El Nino Cycle
(soi.ar = ar(soi, order.max=30)) # estimates and AICs
dev.new()
plot(1:30, soi.ar$aic[-1], type="o") # plot AICs
```

R works only with the AIC in this case. To generate Figure 4.16 we used the following code to obtain AIC, AICc, and BIC. Because AIC and AICc are nearly identical in this example, we only graphed AIC and BIC+1; we added 1 to the BIC to reduce white space in the graphic.

```
n = length(soi)
AIC = rep(0, 30) -> AICc -> BIC
for (k in 1:30){
    sigma2 = ar(soi, order=k, aic=FALSE)$var.pred
    BIC[k] = log(sigma2) + (k*log(n)/n)
    AICc[k] = log(sigma2) + ((n+k)/(n-k-2))
    AIC[k] = log(sigma2) + ((n+2*k)/n)
}
```

Series: soi AR (15) spectrum


Fig. 4.17. Autoregressive spectral estimator for the SOI series using the AR(15) model selected by AIC, AICc, and BIC. The first peak (marked by a vertical dotted line) corresponds to the El Niño period of 52 months.

```
IC = cbind(AIC, BIC+1)
ts.plot(IC, type="o", xlab="p", ylab="AIC / BIC")
```


### 4.7 Linear Filters

Some of the examples of the previous sections have hinted at the possibility the distribution of power or variance in a time series can be modified by making a linear transformation. In this section, we explore that notion further by defining a linear filter and showing how it can be used to extract signals from a time series. The linear filter modifies the spectral characteristics of a time series in a predictable way, and the systematic development of methods for taking advantage of the special properties of linear filters is an important topic in time series analysis.

A linear filter uses a set of specified coefficients $a_{j}$, for $j=0, \pm 1, \pm 2, \ldots$, to transform an input series, $x_{t}$, producing an output series, $y_{t}$, of the form

$$
\begin{equation*}
y_{t}=\sum_{j=-\infty}^{\infty} a_{j} x_{t-j}, \quad \sum_{j=-\infty}^{\infty}\left|a_{j}\right|<\infty . \tag{4.56}
\end{equation*}
$$

The form (4.56) is also called a convolution in some statistical contexts. The coefficients, collectively called the impulse response function, are required to satisfy absolute summability so $y_{t}$ in (4.56) exists as a limit in mean square and the infinite Fourier transform

$$
\begin{equation*}
A_{y x}(\omega)=\sum_{j=-\infty}^{\infty} a_{j} \mathrm{e}^{-2 \pi i \omega j} \tag{4.57}
\end{equation*}
$$

called the frequency response function, is well defined. We have already encountered several linear filters, for example, the simple three-point moving average in Example 1.7, which can be put into the form of (4.56) by letting $a_{-1}=a_{0}=a_{1}=1 / 3$ and taking $a_{t}=0$ for $|j| \geq 2$.

The importance of the linear filter stems from its ability to enhance certain parts of the spectrum of the input series. We now state the following result.

## Property 4.5 Output Spectrum of a Filtered Stationary Series

Assuming existence of spectra, the spectrum of the filtered output $y_{t}$ in (4.56) is related to the spectrum of the input $x_{t}$ by

$$
\begin{equation*}
f_{y y}(\omega)=\left|A_{y x}(\omega)\right|^{2} f_{x x}(\omega) \tag{4.58}
\end{equation*}
$$

where the frequency response function $A_{y x}(\omega)$ is defined in (4.57).
The result (4.58) enables us to calculate the exact effect on the spectrum of any given filtering operation. This important property shows the spectrum of the input series is changed by filtering and the effect of the change can be characterized as a frequency-by-frequency multiplication by the squared magnitude of the frequency response function. Again, an obvious analogy to a property of the variance in classical statistics holds, namely, if $x$ is a random variable with variance $\sigma_{x}^{2}$, then $y=a x$ will have variance $\sigma_{y}^{2}=a^{2} \sigma_{x}^{2}$, so the variance of the linearly transformed random variable is changed by multiplication by $a^{2}$ in much the same way as the linearly filtered spectrum is changed in (4.58).

Finally, we mention that Property 4.2, which was used to get the spectrum of an ARMA process, is just a special case of Property 4.5 where in (4.56), $x_{t}=w_{t}$ is white noise, in which case $f_{x x}(\omega)=\sigma_{w}^{2}$, and $a_{j}=\psi_{j}$, in which case

$$
A_{y x}(\omega)=\psi\left(\mathrm{e}^{-2 \pi i \omega}\right)=\theta\left(\mathrm{e}^{-2 \pi i \omega}\right) / \phi\left(\mathrm{e}^{-2 \pi i \omega}\right) .
$$

## Example 4.13 First Difference and Moving Average Filters

We illustrate the effect of filtering with two common examples, the first difference filter

$$
y_{t}=\nabla x_{t}=x_{t}-x_{t-1}
$$

and the symmetric moving average filter

$$
y_{t}=\frac{1}{24}\left(x_{t-6}+x_{t+6}\right)+\frac{1}{12} \sum_{r=-5}^{5} x_{t-r},
$$

which is a modified Daniell kernel with $m=6$. The results of filtering the SOI series using the two filters are shown in the middle and bottom panels of Figure 4.18. Notice that the effect of differencing is to roughen the series because it tends to


Fig. 4.18. SOI series (top) compared with the differenced SOI (middle) and a centered 12month moving average (bottom).
retain the higher or faster frequencies. The centered moving average smoothes the series because it retains the lower frequencies and tends to attenuate the higher frequencies. In general, differencing is an example of a high-pass filter because it retains or passes the higher frequencies, whereas the moving average is a low-pass filter because it passes the lower or slower frequencies.

Notice that the slower periods are enhanced in the symmetric moving average and the seasonal or yearly frequencies are attenuated. The filtered series makes about 9 cycles in the length of the data (about one cycle every 52 months) and the moving average filter tends to enhance or extract the signal that is associated with El Niño. Moreover, by the low-pass filtering of the data, we get a better sense of the El Niño effect and its irregularity. Figure 4.19 shows the results of a spectral analysis on the low-pass filtered SOI series. It is clear that all high frequency behavior has been removed and the El Niño cycle is accentuated; the dotted vertical line in the figure corresponds to the 52 months cycle.

Now, having done the filtering, it is essential to determine the exact way in which the filters change the input spectrum. We shall use (4.57) and (4.58) for this purpose. The first difference filter can be written in the form (4.56) by letting $a_{0}=1, a_{1}=-1$, and $a_{r}=0$ otherwise. This implies that

$$
A_{y x}(\omega)=1-\mathrm{e}^{-2 \pi i \omega}
$$

SOI - Twelve Month Moving Average


Fig. 4.19. Spectral analysis of SOI after applying a 12 -month moving average filter. The vertical line corresponds to the 52 -month cycle.
and the squared frequency response becomes

$$
\begin{equation*}
\left|A_{y x}(\omega)\right|^{2}=\left(1-\mathrm{e}^{-2 \pi i \omega}\right)\left(1-\mathrm{e}^{2 \pi i \omega}\right)=2[1-\cos (2 \pi \omega)] . \tag{4.59}
\end{equation*}
$$

The top panel of Figure 4.20 shows that the first difference filter will attenuate the lower frequencies and enhance the higher frequencies because the multiplier of the spectrum, $\left|A_{y x}(\omega)\right|^{2}$, is large for the higher frequencies and small for the lower frequencies. Generally, the slow rise of this kind of filter does not particularly recommend it as a procedure for retaining only the high frequencies.

For the centered 12 -month moving average, we can take $a_{-6}=a_{6}=1 / 24$, $a_{k}=1 / 12$ for $-5 \leq k \leq 5$ and $a_{k}=0$ elsewhere. Substituting and recognizing the cosine terms gives

$$
\begin{equation*}
A_{y x}(\omega)=\frac{1}{12}\left[1+\cos (12 \pi \omega)+2 \sum_{k=1}^{5} \cos (2 \pi \omega k)\right] \tag{4.60}
\end{equation*}
$$

Plotting the squared frequency response of this function as in Figure 4.20 shows that we can expect this filter to cut most of the frequency content above .05 cycles per point. This corresponds to eliminating periods shorter than $T=1 / .05=20$ points. In particular, this drives down the yearly components with periods of $T=12$ months and enhances the El Niño frequency, which is somewhat lower. The filter is not completely efficient at attenuating high frequencies; some power contributions are left at higher frequencies, as shown in the function $\left|A_{y x}(\omega)\right|^{2}$ and in the spectrum of the moving average shown in Figure 4.4.

The following R session shows how to filter the data, perform the spectral analysis of this example, and plot the squared frequency response curve of the difference filter.


Fig. 4.20. Squared frequency response functions of the first difference and 12-month moving average filters.

```
par(mfrow=c(3,1))
plot(soi) # plot data
plot(diff(soi)) # plot first difference
k = kernel("modified.daniell", 6) # filter weights
plot(soif <- kernapply(soi, k)) # plot 12 month filter
dev.new()
mvspec(soif, spans=9, taper=.1, log="no") # spectral analysis
abline(v=12/52, lty="dashed")
dev.new()
w = seq(0, .5, length=500) # frequency response
FR = abs(1-exp(2i*pi*w))^2
plot(w, FR, type="l")
```


### 4.8 Multiple Series and Cross-Spectra

The notion of analyzing frequency fluctuations using classical statistical ideas extends to the case in which there are several jointly stationary series, for example, $x_{t}$ and $y_{t}$. In this case, we can introduce the idea of a correlation indexed by frequency, called the coherence. The autocovariance function

$$
\gamma_{x y}(h)=E\left[\left(x_{t+h}-\mu_{x}\right)\left(y_{t}-\mu_{y}\right)\right]
$$

has a spectral representation given by

$$
\begin{equation*}
\gamma_{x y}(h)=\int_{-1 / 2}^{1 / 2} f_{x y}(\omega) \mathrm{e}^{2 \pi i \omega h} d \omega \quad h=0, \pm 1, \pm 2, \ldots \tag{4.61}
\end{equation*}
$$

where the cross-spectrum is defined as the Fourier transform

$$
\begin{equation*}
f_{x y}(\omega)=\sum_{h=-\infty}^{\infty} \gamma_{x y}(h) \mathrm{e}^{-2 \pi i \omega h} \quad-1 / 2 \leq \omega \leq 1 / 2 \tag{4.62}
\end{equation*}
$$

assuming that the cross-covariance function is absolutely summable, as was the case for the autocovariance. The cross-spectrum is generally a complex-valued function, and it is often written as ${ }^{5}$

$$
\begin{equation*}
f_{x y}(\omega)=c_{x y}(\omega)-i q_{x y}(\omega) \tag{4.63}
\end{equation*}
$$

where

$$
\begin{equation*}
c_{x y}(\omega)=\sum_{h=-\infty}^{\infty} \gamma_{x y}(h) \cos (2 \pi \omega h) \tag{4.64}
\end{equation*}
$$

and

$$
\begin{equation*}
q_{x y}(\omega)=\sum_{h=-\infty}^{\infty} \gamma_{x y}(h) \sin (2 \pi \omega h) \tag{4.65}
\end{equation*}
$$

are defined as the cospectrum and quadspectrum, respectively. Because of the relationship $\gamma_{y x}(h)=\gamma_{x y}(-h)$, it follows, by substituting into (4.62) and rearranging, that

$$
\begin{equation*}
f_{y x}(\omega)=\overline{f_{x y}(\omega)} \tag{4.66}
\end{equation*}
$$

This result, in turn, implies that the cospectrum and quadspectrum satisfy

$$
\begin{equation*}
c_{y x}(\omega)=c_{x y}(\omega) \tag{4.67}
\end{equation*}
$$

and

$$
\begin{equation*}
q_{y x}(\omega)=-q_{x y}(\omega) . \tag{4.68}
\end{equation*}
$$

An important example of the application of the cross-spectrum is to the problem of predicting an output series $y_{t}$ from some input series $x_{t}$ through a linear filter relation such as the three-point moving average considered below. A measure of the strength of such a relation is the squared coherence function, defined as

$$
\begin{equation*}
\rho_{y \cdot x}^{2}(\omega)=\frac{\left|f_{y x}(\omega)\right|^{2}}{f_{x x}(\omega) f_{y y}(\omega)} \tag{4.69}
\end{equation*}
$$

where $f_{x x}(\omega)$ and $f_{y y}(\omega)$ are the individual spectra of the $x_{t}$ and $y_{t}$ series, respectively. Although we consider a more general form of this that applies to multiple inputs later, it is instructive to display the single input case as (4.69) to emphasize the analogy with conventional squared correlation, which takes the form

$$
\rho_{y x}^{2}=\frac{\sigma_{y x}^{2}}{\sigma_{x}^{2} \sigma_{y}^{2}},
$$

for random variables with variances $\sigma_{x}^{2}$ and $\sigma_{y}^{2}$ and covariance $\sigma_{y x}=\sigma_{x y}$. This motivates the interpretation of squared coherence and the squared correlation between two time series at frequency $\omega$.

[^23]
## Example 4.14 Three-Point Moving Average

As a simple example, we compute the cross-spectrum between $x_{t}$ and the threepoint moving average $y_{t}=\left(x_{t-1}+x_{t}+x_{t+1}\right) / 3$, where $x_{t}$ is a stationary input process with spectral density $f_{x x}(\omega)$. First,

$$
\begin{aligned}
\gamma_{x y}(h) & =\operatorname{cov}\left(x_{t+h}, y_{t}\right)=\frac{1}{3} \operatorname{cov}\left(x_{t+h}, x_{t-1}+x_{t}+x_{t+1}\right) \\
& =\frac{1}{3}\left(\gamma_{x x}(h+1)+\gamma_{x x}(h)+\gamma_{x x}(h-1)\right) \\
& =\frac{1}{3} \int_{-1 / 2}^{1 / 2}\left(\mathrm{e}^{2 \pi i \omega}+1+\mathrm{e}^{-2 \pi i \omega}\right) \mathrm{e}^{2 \pi i \omega h} f_{x x}(\omega) d \omega \\
& =\frac{1}{3} \int_{-1 / 2}^{1 / 2}[1+2 \cos (2 \pi \omega)] f_{x x}(\omega) \mathrm{e}^{2 \pi i \omega h} d \omega
\end{aligned}
$$

where we have use (4.13). Using the uniqueness of the Fourier transform, we argue from the spectral representation (4.61) that

$$
f_{x y}(\omega)=\frac{1}{3}[1+2 \cos (2 \pi \omega)] f_{x x}(\omega)
$$

so that the cross-spectrum is real in this case. From Example 4.5, the spectral density of $y_{t}$ is

$$
\begin{aligned}
f_{y y}(\omega) & =\frac{1}{9}[3+4 \cos (2 \pi \omega)+2 \cos (4 \pi \omega)] f_{x x}(\omega) \\
& =\frac{1}{9}[1+2 \cos (2 \pi \omega)]^{2} f_{x x}(\omega),
\end{aligned}
$$

using the identity $\cos (2 \alpha)=2 \cos ^{2}(\alpha)-1$ in the last step. Substituting into (4.69) yields the squared coherence between $x_{t}$ and $y_{t}$ as unity over all frequencies. This is a characteristic inherited by more general linear filters. However, if some noise is added to the three-point moving average, the coherence is not unity; these kinds of models will be considered in detail later.

## Property 4.6 Spectral Representation of a Vector Process

If the elements of the $p \times p$ autocovariance function matrix

$$
\Gamma(h)=E\left[\left(\boldsymbol{x}_{t+h}-\boldsymbol{\mu}\right)\left(\boldsymbol{x}_{t}-\boldsymbol{\mu}\right)^{\prime}\right]
$$

of a p-dimensional stationary time series, $\boldsymbol{x}_{t}=\left(x_{t 1}, x_{t 2}, \ldots, x_{t p}\right)^{\prime}$, has elements satisfying

$$
\begin{equation*}
\sum_{h=-\infty}^{\infty}\left|\gamma_{j k}(h)\right|<\infty \tag{4.70}
\end{equation*}
$$

for all $j, k=1, \ldots, p$, then $\Gamma(h)$ has the representation

$$
\begin{equation*}
\Gamma(h)=\int_{-1 / 2}^{1 / 2} \mathrm{e}^{2 \pi i \omega h} f(\omega) d \omega \quad h=0, \pm 1, \pm 2, \ldots \tag{4.71}
\end{equation*}
$$

as the inverse transform of the spectral density matrix, $f(\omega)=\left\{f_{j k}(\omega)\right\}$, for $j, k=$ $1, \ldots, p$, with elements equal to the cross-spectral components. The matrix $f(\omega)$ has the representation

$$
\begin{equation*}
f(\omega)=\sum_{h=-\infty}^{\infty} \Gamma(h) \mathrm{e}^{-2 \pi i \omega h} \quad-1 / 2 \leq \omega \leq 1 / 2 \tag{4.72}
\end{equation*}
$$

## Example 4.15 Spectral Matrix of a Bivariate Process

Consider a jointly stationary bivariate process $\left(x_{t}, y_{t}\right)$. We arrange the autocovariances in the matrix

$$
\Gamma(h)=\left(\begin{array}{ll}
\gamma_{x x}(h) & \gamma_{x y}(h) \\
\gamma_{y x}(h) & \gamma_{y y}(h)
\end{array}\right) .
$$

The spectral matrix would be given by

$$
f(\omega)=\left(\begin{array}{ll}
f_{x x}(\omega) & f_{x y}(\omega) \\
f_{y x}(\omega) & f_{y y}(\omega)
\end{array}\right)
$$

where the Fourier transform (4.71) and (4.72) relate the autocovariance and spectral matrices.

The extension of spectral estimation to vector series is fairly obvious. For the vector series $\boldsymbol{x}_{t}=\left(x_{t 1}, x_{t 2}, \ldots, x_{t p}\right)^{\prime}$, we may use the vector of DFTs, say $\boldsymbol{d}\left(\omega_{j}\right)=$ $\left(d_{1}\left(\omega_{j}\right), d_{2}\left(\omega_{j}\right), \ldots, d_{p}\left(\omega_{j}\right)\right)^{\prime}$, and estimate the spectral matrix by

$$
\begin{equation*}
\bar{f}(\omega)=L^{-1} \sum_{k=-m}^{m} I\left(\omega_{j}+k / n\right) \tag{4.73}
\end{equation*}
$$

where now

$$
\begin{equation*}
I\left(\omega_{j}\right)=\boldsymbol{d}\left(\omega_{j}\right) \boldsymbol{d}^{*}\left(\omega_{j}\right) \tag{4.74}
\end{equation*}
$$

is a $p \times p$ complex matrix. ${ }^{6}$
Again, the series may be tapered before the DFT is taken in (4.73) and we can use weighted estimation,

$$
\begin{equation*}
\widehat{f}(\omega)=\sum_{k=-m}^{m} h_{k} I\left(\omega_{j}+k / n\right) \tag{4.75}
\end{equation*}
$$

where $\left\{h_{k}\right\}$ are weights as defined in (4.42). The estimate of squared coherence between two series, $y_{t}$ and $x_{t}$ is

$$
\begin{equation*}
\widehat{\rho}_{y \cdot x}^{2}(\omega)=\frac{\left|\widehat{f}_{y x}(\omega)\right|^{2}}{\widehat{f}_{x x}(\omega) \widehat{f}_{y y}(\omega)} \tag{4.76}
\end{equation*}
$$

If the spectral estimates in (4.76) are obtained using equal weights, we will write $\bar{\rho}_{y \cdot x}^{2}(\omega)$ for the estimate.

Under general conditions, if $\rho_{y \cdot x}^{2}(\omega)>0$ then

[^24]

Fig. 4.21. Squared coherency between the SOI and Recruitment series; $L=19, n=453, n^{\prime}=$ 480 , and $\alpha=.001$. The horizontal line is $C .001$.

$$
\begin{equation*}
\left|\widehat{\rho}_{y \cdot x}(\omega)\right| \sim A N\left(\left|\rho_{y \cdot x}(\omega)\right|,\left(1-\rho_{y \cdot x}^{2}(\omega)\right)^{2} / 2 L_{h}\right) \tag{4.77}
\end{equation*}
$$

where $L_{h}$ is defined in (4.44); the details of this result may be found in Brockwell and Davis (1991, Ch 11). We may use (4.77) to obtain approximate confidence intervals for the squared coherency $\rho_{y \cdot x}^{2}(\omega)$.

We can test the hypothesis that $\rho_{y \cdot x}^{2}(\omega)=0$ if we use $\bar{\rho}_{y \cdot x}^{2}(\omega)$ for the estimate with $L>1,{ }^{7}$ that is,

$$
\begin{equation*}
\bar{\rho}_{y \cdot x}^{2}(\omega)=\frac{\left|\bar{f}_{y x}(\omega)\right|^{2}}{\bar{f}_{x x}(\omega) \bar{f}_{y y}(\omega)} . \tag{4.78}
\end{equation*}
$$

In this case, under the null hypothesis, the statistic

$$
\begin{equation*}
F=\frac{\bar{\rho}_{y \cdot x}^{2}(\omega)}{\left(1-\bar{\rho}_{y \cdot x}^{2}(\omega)\right)}(L-1) \tag{4.79}
\end{equation*}
$$

has an approximate $F$-distribution with 2 and $2 L-2$ degrees of freedom. When the series have been extended to length $n^{\prime}$, we replace $2 L-2$ by $d f-2$, where $d f$ is defined in (4.39). Solving (4.79) for a particular significance level $\alpha$ leads to

$$
\begin{equation*}
C_{\alpha}=\frac{F_{2,2 L-2}(\alpha)}{L-1+F_{2,2 L-2}(\alpha)} \tag{4.80}
\end{equation*}
$$

as the approximate value that must be exceeded for the original squared coherence to be able to reject $\rho_{y \cdot x}^{2}(\omega)=0$ at an a priori specified frequency.

## Example 4.16 Coherence Between SOI and Recruitment

Figure 4.21 shows the squared coherence between the SOI and Recruitment series over a wider band than was used for the spectrum. In this case, we used $L=19, d f=$

[^25]$2(19)(453 / 480) \approx 36$ and $F_{2, d f-2}(.001) \approx 8.53$ at the significance level $\alpha=.001$. Hence, we may reject the hypothesis of no coherence for values of $\bar{\rho}_{y \cdot x}^{2}(\omega)$ that exceed $C_{.001}=.32$. We emphasize that this method is crude because, in addition to the fact that the $F$-statistic is approximate, we are examining the squared coherence across all frequencies with the Bonferroni inequality in mind. Figure 4.21 also exhibits confidence bands as part of the R plotting routine. We emphasize that these bands are only valid for $\omega$ where $\rho_{y \cdot x}^{2}(\omega)>0$.

In this case, the seasonal frequency and the El Niño frequencies ranging between about 3 and 7 year periods are strongly coherent. Other frequencies are also strongly coherent, although the strong coherence is less impressive because the underlying power spectrum at these higher frequencies is fairly small. Finally, we note that the coherence is persistent at the seasonal harmonic frequencies.

This example may be reproduced using the following R commands.
sr = mvspec(cbind(soi,rec), kernel('daniell', 9), plot=FALSE)
sr\$df \# = 35.8625
$\mathrm{f}=\mathrm{qf}(.999,2, \mathrm{sr} \$ \mathrm{df}-2) \quad \#=8.529792$
$\mathrm{C}=\mathrm{f} /(18+\mathrm{f}) \quad \#=0.318878$
plot(sr, plot.type = "coh", ci.lty = 2)
abline (h = C)

## Problems

4.1 Repeat the simulations and analyses in Example 4.1 and Example 4.2 with the following changes:
(a) Change the sample size to $n=128$ and generate and plot the same series as in Example 4.1:

$$
\begin{aligned}
x_{t 1} & =2 \cos (2 \pi .06 t)+3 \sin (2 \pi .06 t), \\
x_{t 2} & =4 \cos (2 \pi .10 t)+5 \sin (2 \pi .10 t), \\
x_{t 3} & =6 \cos (2 \pi .40 t)+7 \sin (2 \pi .40 t), \\
x_{t} & =x_{t 1}+x_{t 2}+x_{t 3} .
\end{aligned}
$$

What is the major difference between these series and the series generated in Example 4.1? (Hint: The answer is fundamental. But if your answer is the series are longer, you may be punished severely.)
(b) As in Example 4.2, compute and plot the periodogram of the series, $x_{t}$, generated in (a) and comment.
(c) Repeat the analyses of (a) and (b) but with $n=100$ (as in Example 4.1), and adding noise to $x_{t}$; that is

$$
x_{t}=x_{t 1}+x_{t 2}+x_{t 3}+w_{t}
$$

where $w_{t} \sim \operatorname{iid} \mathrm{~N}\left(0, \sigma_{w}=5\right)$. That is, you should simulate and plot the data, and then plot the periodogram of $x_{t}$ and comment.
4.2 Verify (4.5).
4.3 Consider an MA(1) process

$$
x_{t}=w_{t}+\theta w_{t-1},
$$

where $\theta$ is a parameter.
(a) Derive a formula for the power spectrum of $x_{t}$, expressed in terms of $\theta$ and $\omega$.
(b) Use arma.spec() to plot the spectral density of $x_{t}$ for $\theta>0$ and for $\theta<0$ (just select arbitrary values).
(c) How should we interpret the spectra exhibited in part (b)?
4.4 Consider a first-order autoregressive model

$$
x_{t}=\phi x_{t-1}+w_{t},
$$

where $\phi$, for $|\phi|<1$, is a parameter and the $w_{t}$ are independent random variables with mean zero and variance $\sigma_{w}^{2}$.
(a) Show that the power spectrum of $x_{t}$ is given by

$$
f_{x}(\omega)=\frac{\sigma_{w}^{2}}{1+\phi^{2}-2 \phi \cos (2 \pi \omega)} .
$$

(b) Verify the autocovariance function of this process is

$$
\gamma_{x}(h)=\frac{\sigma_{w}^{2} \phi^{|h|}}{1-\phi^{2}}
$$

$h=0, \pm 1, \pm 2, \ldots$, by showing that the inverse transform of $\gamma_{x}(h)$ is the spectrum derived in part (a).
4.5 In applications, we will often observe series containing a signal that has been delayed by some unknown time $D$, i.e.,

$$
x_{t}=s_{t}+A s_{t-D}+n_{t},
$$

where $s_{t}$ and $n_{t}$ are stationary and independent with zero means and spectral densities $f_{s}(\omega)$ and $f_{n}(\omega)$, respectively. The delayed signal is multiplied by some unknown constant $A$. Find the autocovariance function of $x_{t}$ and use it to show

$$
f_{x}(\omega)=\left[1+A^{2}+2 A \cos (2 \pi \omega D)\right] f_{s}(\omega)+f_{n}(\omega) .
$$

4.6 Figure 4.22 shows the biyearly smoothed (12-month moving average) number of sunspots from June 1749 to December 1978 with $n=459$ points that were taken twice per year; the data are contained in sunspotz. With Example 4.7 as a guide, perform a periodogram analysis identifying the predominant periods and obtaining confidence intervals for the identified periods. Interpret your findings.


Fig. 4.22. Smoothed 12 -month sunspot numbers (sunspotz) sampled twice per year.
4.7 The levels of salt concentration known to have occurred over rows, corresponding to the average temperature levels for the soil science are in salt and saltemp. Plot the series and then identify the dominant frequencies by performing separate spectral analyses on the two series. Include confidence intervals for the dominant frequencies and interpret your findings.
4.8 Let the observed series $x_{t}$ be composed of a periodic signal and noise so it can be written as

$$
x_{t}=\beta_{1} \cos \left(2 \pi \omega_{k} t\right)+\beta_{2} \sin \left(2 \pi \omega_{k} t\right)+w_{t},
$$

where $w_{t}$ is a white noise process with variance $\sigma_{w}^{2}$. The frequency $\omega_{k} \neq 0, \frac{1}{2}$ is assumed to be known and of the form $k / n$. Given data $x_{1}, \ldots, x_{n}$, suppose we consider estimating $\beta_{1}, \beta_{2}$ and $\sigma_{w}^{2}$ by least squares.
(a) Use simple regression formulas to show that for a fixed $\omega_{k}$, the least squares regression coefficients are

$$
\widehat{\beta}_{1}=2 n^{-1 / 2} d_{c}\left(\omega_{k}\right) \quad \text { and } \quad \widehat{\beta}_{2}=2 n^{-1 / 2} d_{s}\left(\omega_{k}\right),
$$

where the cosine and sine transforms (4.21) and (4.22) appear on the right-hand side. Hint: See Problem 4.20.
(b) Prove that the error sum of squares can be written as

$$
S S E=\sum_{t=1}^{n} x_{t}^{2}-2 I_{x}\left(\omega_{k}\right)
$$

so that the value of $\omega_{k}$ that minimizes squared error is the same as the value that maximizes the periodogram $I_{x}\left(\omega_{k}\right)$ estimator (4.19).
(c) Show that the sum of squares for the regression is given by

$$
S S R=2 I_{x}\left(\omega_{k}\right)
$$

(d) Under the Gaussian assumption and fixed $\omega_{k}$, show that the $F$-test of no regression leads to an $F$-statistic that is a monotone function of $I_{x}\left(\omega_{k}\right)$.
4.9 Repeat Problem 4.6 using a nonparametric spectral estimation procedure. In addition to discussing your findings in detail, comment on your choice of a spectral estimate with regard to smoothing and tapering.
4.10 Repeat Problem 4.7 using a nonparametric spectral estimation procedure. In addition to discussing your findings in detail, comment on your choice of a spectral estimate with regard to smoothing and tapering.
4.11 Often, the periodicities in the sunspot series are investigated by fitting an autoregressive spectrum of sufficiently high order. The main periodicity is often stated to be in the neighborhood of 11 years. Fit an autoregressive spectral estimator to the sunspot data using a model selection method of your choice. Compare the result with a conventional nonparametric spectral estimator found in Problem 4.6.
4.12 Fit an autoregressive spectral estimator to the Recruitment series and compare it to the results of Example 4.10.
4.13 The periodic behavior of a time series induced by echoes can also be observed in the spectrum of the series; this fact can be seen from the results stated in Problem 4.5(a). Using the notation of that problem, suppose we observe $x_{t}=s_{t}+A s_{t-D}+$ $n_{t}$, which implies the spectra satisfy $f_{x}(\omega)=\left[1+A^{2}+2 A \cos (2 \pi \omega D)\right] f_{s}(\omega)+f_{n}(\omega)$. If the noise is negligible $\left(f_{n}(\omega) \approx 0\right)$ then $\log f_{x}(\omega)$ is approximately the sum of a periodic component, $\log \left[1+A^{2}+2 A \cos (2 \pi \omega D)\right]$, and $\log f_{s}(\omega)$. Bogart et al. (1962) proposed treating the detrended $\log$ spectrum as a pseudo time series and calculating its spectrum, or cepstrum, which should show a peak at a quefrency corresponding to $1 / D$. The cepstrum can be plotted as a function of quefrency, from which the delaty $D$ can be estimated.

For the speech series presented in speech, estimate the pitch period using cepstral analysis as follows.
(a) Calculate and display the log-periodogram of the data. Is the periodogram periodic, as predicted?
(b) Perform a cepstral (spectral) analysis on the detrended logged periodogram, and use the results to estimate the delay $D$. How does your answer compare with the analysis of ??, which was based on the ACF?
4.14 Consider two time series

$$
\begin{gathered}
x_{t}=w_{t}-w_{t-1}, \\
y_{t}=\frac{1}{2}\left(w_{t}+w_{t-1}\right),
\end{gathered}
$$

formed from the white noise series $w_{t}$ with variance $\sigma_{w}^{2}=1$.
(a) Are $x_{t}$ and $y_{t}$ jointly stationary? Recall the cross-covariance function must also be a function only of the lag $h$ and cannot depend on time.
(b) Compute the spectra $f_{y}(\omega)$ and $f_{x}(\omega)$, and comment on the difference between the two results.
(c) Suppose sample spectral estimators $\bar{f}_{y}(.10)$ are computed for the series using $L=3$. Find $a$ and $b$ such that

$$
P\left\{a \leq \bar{f}_{y}(.10) \leq b\right\}=.90
$$

This expression gives two points that will contain $90 \%$ of the sample spectral values. Put $5 \%$ of the area in each tail.
4.15 Analyze the coherency between the temperature and salt data discussed in Problem 4.7. Discuss your findings.

### 4.16 Consider two processes

$$
x_{t}=w_{t} \quad \text { and } \quad y_{t}=\phi x_{t-D}+v_{t}
$$

where $w_{t}$ and $v_{t}$ are independent white noise processes with common variance $\sigma^{2}, \phi$ is a constant, and $D$ is a fixed integer delay.
(a) Compute the coherency between $x_{t}$ and $y_{t}$.
(b) Simulate $n=1024$ normal observations from $x_{t}$ and $y_{t}$ for $\phi=.9, \sigma^{2}=1$, and $D=0$. Then estimate and plot the coherency between the simulated series for the following values of $L$ and comment:
(i) $L=1$, (ii) $L=3$, (iii) $L=41$, and (iv) $L=101$.
4.17 For the processes in Problem 4.16:
(a) Compute the phase between $x_{t}$ and $y_{t}$.
(b) Simulate $n=1024$ observations from $x_{t}$ and $y_{t}$ for $\phi=.9, \sigma^{2}=1$, and $D=1$. Then estimate and plot the phase between the simulated series for the following values of $L$ and comment:
(i) $L=1$, (ii) $L=3$, (iii) $L=41$, and (iv) $L=101$.
4.18 Consider the bivariate time series records containing monthly U.S. production as measured by the Federal Reserve Board Production Index (prodn) and monthly unemployment (unemp) that are included with astsa.
(a) Compute the spectrum and the log spectrum for each series, and identify statistically significant peaks. Explain what might be generating the peaks. Compute the coherence, and explain what is meant when a high coherence is observed at a particular frequency.
(b) What would be the effect of applying the filter

$$
u_{t}=x_{t}-x_{t-1} \quad \text { followed by } \quad v_{t}=u_{t}-u_{t-12}
$$

to the series given above? Plot the predicted frequency responses of the simple difference filter and of the seasonal difference of the first difference.
(c) Apply the filters successively to one of the two series and plot the output. Examine the output after taking a first difference and comment on whether stationarity is a reasonable assumption. Why or why not? Plot after taking the seasonal difference of the first difference. What can be noticed about the output that is consistent with what you have predicted from the frequency response? Verify by computing the spectrum of the output after filtering.
4.19 Let $x_{t}=\cos (2 \pi \omega t)$, and consider the output $y_{t}=\sum_{k=-\infty}^{\infty} a_{k} x_{t-k}$, where $\sum_{k}\left|a_{k}\right|<\infty$. Show $y_{t}=|A(\omega)| \cos (2 \pi \omega t+\phi(\omega))$, where $|A(\omega)|$ and $\phi(\omega)$ are the amplitude and phase of the filter, respectively. Interpret the result in terms of the relationship between the input series, $x_{t}$, and the output series, $y_{t}$.
4.20 * This is here for useful information. Verify that for any positive integer $n$ and $j, k=0,1, \ldots, \llbracket n / 2 \rrbracket$, where $\llbracket \cdot \rrbracket$ denotes the greatest integer function:
(a) Except for $j=0$ or $j=n / 2$,

$$
\sum_{t=1}^{n} \cos ^{2}(2 \pi t j / n)=\sum_{t=1}^{n} \sin ^{2}(2 \pi t j / n)=n / 2 .
$$

(b) When $j=0$ or $j=n / 2$,

$$
\sum_{t=1}^{n} \cos ^{2}(2 \pi t j / n)=n \text { but } \sum_{t=1}^{n} \sin ^{2}(2 \pi t j / n)=0 .
$$

(c) For $j \neq k$,

$$
\sum_{t=1}^{n} \cos (2 \pi t j / n) \cos (2 \pi t k / n)=\sum_{t=1}^{n} \sin (2 \pi t j / n) \sin (2 \pi t k / n)=0 .
$$

(d) Also, for any $j$ and $k$,

$$
\sum_{t=1}^{n} \cos (2 \pi t j / n) \sin (2 \pi t k / n)=0
$$

$$
\begin{aligned}
& \text { * Note, } \sum_{t=1}^{n} z^{t}=z \frac{1-z^{n}}{1-z} \text { for } z \neq 1 \text {, and we'll do (a): } \\
& \begin{aligned}
\sum_{t=1}^{n} \cos ^{2}(2 \pi t j / n) & =\frac{1}{4} \sum_{t=1}^{n}\left(\mathrm{e}^{2 \pi i t j / n}+n+\mathrm{e}^{-2 \pi i t j / n}\right)\left(\mathrm{e}^{2 \pi i t j / n}+\mathrm{e}^{-2 \pi i t j / n}\right) \\
& =\frac{1}{4} \sum_{t=1}^{n}\left(\mathrm{e}^{4 \pi i t j / n}+1+1+\mathrm{e}^{-4 \pi i t j / n}\right)=\frac{n}{2} .
\end{aligned}
\end{aligned}
$$

## Chapter 5

## Some Additional Topics **

In this chapter, we present special or advanced topics in the time domain. This chapter consists of sections of independent topics that may be read in any order. Most of the sections depend on a basic knowledge of ARMA models, forecasting and estimation, which is the material that is covered in Chapter 3. A few sections, for example the section on long memory models, require some knowledge of spectral analysis and related topics covered in Chapter 4.

### 5.1 GARCH Models

Recent problems in finance have motivated the study of the volatility, or variability, of a time series. Although ARMA models assume a constant variance, models such as the autoregressive conditionally heteroscedastic or ARCH model, first introduced by Engle (1982), were developed to model changes in volatility. These models were later extended to generalized ARCH, or GARCH models by Bollerslev (1986).

In $\S 3.8$, we discussed the return or growth rate of a series. For example, if $x_{t}$ is the value of a stock at time $t$, then the return or relative gain, $y_{t}$, of the stock at time $t$ is

$$
\begin{equation*}
y_{t}=\frac{x_{t}-x_{t-1}}{x_{t-1}} . \tag{5.1}
\end{equation*}
$$

Definition (5.1) implies that $x_{t}=\left(1+y_{t}\right) x_{t-1}$. Thus, based on the discussion in $\S 3.8$, if the return represents a small (in magnitude) percentage change then

$$
\begin{equation*}
\nabla\left[\log \left(x_{t}\right)\right] \approx y_{t} . \tag{5.2}
\end{equation*}
$$

Either value, $\nabla\left[\log \left(x_{t}\right)\right]$ or $\left(x_{t}-x_{t-1}\right) / x_{t-1}$, will be called the return, and will be denoted by $y_{t}$. It is the study of $y_{t}$ that is the focus of ARCH, GARCH, and other volatility models. Recently there has been interest in stochastic volatility models and we will discuss these models in Chapter 6 because they are state-space models.

Typically, for financial series, the return $y_{t}$, does not have a constant conditional variance, and highly volatile periods tend to be clustered together. In other words,
there is a strong dependence of sudden bursts of variability in a return on the series own past. For example, Figure 1.4 shows the daily returns of the New York Stock Exchange (NYSE) from February 2, 1984 to December 31, 1991. In this case, as is typical, the return $y_{t}$ is fairly stable, except for short-term bursts of high volatility.

The simplest ARCH model, the $\operatorname{ARCH}(1)$, models the return as

$$
\begin{align*}
y_{t} & =\sigma_{t} \epsilon_{t}  \tag{5.3}\\
\sigma_{t}^{2} & =\alpha_{0}+\alpha_{1} y_{t-1}^{2}, \tag{5.4}
\end{align*}
$$

where $\epsilon_{t}$ is standard Gaussian white noise; that is, $\epsilon_{t} \sim \operatorname{iid} \mathrm{~N}(0,1)$. As with ARMA models, we must impose some constraints on the model parameters to obtain desirable properties. One obvious constraint is that $\alpha_{1}$ must not be negative, or else $\sigma_{t}^{2}$ may be negative.

As we shall see, the $\mathrm{ARCH}(1)$ models return as a white noise process with nonconstant conditional variance, and that conditional variance depends on the previous return. First, notice that the conditional distribution of $y_{t}$ given $y_{t-1}$ is Gaussian:

$$
\begin{equation*}
y_{t} \mid y_{t-1} \sim \mathrm{~N}\left(0, \alpha_{0}+\alpha_{1} y_{t-1}^{2}\right) . \tag{5.5}
\end{equation*}
$$

In addition, it is possible to write the $\operatorname{ARCH}(1)$ model as a non-Gaussian $\operatorname{AR}(1)$ model in the square of the returns $y_{t}^{2}$. First, rewrite (5.3)-(5.4) as

$$
\begin{aligned}
y_{t}^{2} & =\sigma_{t}^{2} \epsilon_{t}^{2} \\
\alpha_{0}+\alpha_{1} y_{t-1}^{2} & =\sigma_{t}^{2},
\end{aligned}
$$

and subtract the two equations to obtain

$$
y_{t}^{2}-\left(\alpha_{0}+\alpha_{1} y_{t-1}^{2}\right)=\sigma_{t}^{2} \epsilon_{t}^{2}-\sigma_{t}^{2}
$$

Now, write this equation as

$$
\begin{equation*}
y_{t}^{2}=\alpha_{0}+\alpha_{1} y_{t-1}^{2}+v_{t} \tag{5.6}
\end{equation*}
$$

where $v_{t}=\sigma_{t}^{2}\left(\epsilon_{t}^{2}-1\right)$. Because $\epsilon_{t}^{2}$ is the square of a $\mathrm{N}(0,1)$ random variable, $\epsilon_{t}^{2}-1$ is a shifted (to have mean-zero), $\chi_{1}^{2}$ random variable.

To explore the properties of ARCH, we define $\mathcal{Y}_{s}=\left\{y_{s}, y_{s-1}, \ldots\right\}$. Then, using (5.5), we immediately see that $y_{t}$ has a zero mean:

$$
\begin{equation*}
E\left(y_{t}\right)=E E\left(y_{t} \mid \mathcal{Y}_{t-1}\right)=E E\left(y_{t} \mid y_{t-1}\right)=0 . \tag{5.7}
\end{equation*}
$$

Because $E\left(y_{t} \mid Y_{t-1}\right)=0$, the process $y_{t}$ is said to be a martingale difference.
Because $y_{t}$ is a martingale difference, it is also an uncorrelated sequence. For example, with $h>0$,

$$
\begin{align*}
\operatorname{cov}\left(y_{t+h}, y_{t}\right) & =E\left(y_{t} y_{t+h}\right)=E E\left(y_{t} y_{t+h} \mid \mathcal{Y}_{t+h-1}\right) \\
& =E\left\{y_{t} E\left(y_{t+h} \mid \mathcal{Y}_{t+h-1}\right)\right\}=0 \tag{5.8}
\end{align*}
$$

The last line of (5.8) follows because $y_{t}$ belongs to the information set $\mathcal{Y}_{t+h-1}$ for $h>0$, and, $E\left(y_{t+h} \mid \mathcal{Y}_{t+h-1}\right)=0$, as determined in (5.7).

An argument similar to (5.7) and (5.8) will establish the fact that the error process $v_{t}$ in (5.6) is also a martingale difference and, consequently, an uncorrelated sequence. If the variance of $v_{t}$ is finite and constant with respect to time, and $0 \leq \alpha_{1}<1$, then based on Property 3.1, (5.6) specifies a causal AR(1) process for $y_{t}^{2}$. Therefore, $E\left(y_{t}^{2}\right)$ and $\operatorname{var}\left(y_{t}^{2}\right)$ must be constant with respect to time $t$. This, implies that

$$
\begin{equation*}
E\left(y_{t}^{2}\right)=\operatorname{var}\left(y_{t}\right)=\frac{\alpha_{0}}{1-\alpha_{1}} \tag{5.9}
\end{equation*}
$$

and, after some manipulations,

$$
\begin{equation*}
E\left(y_{t}^{4}\right)=\frac{3 \alpha_{0}^{2}}{\left(1-\alpha_{1}\right)^{2}} \frac{1-\alpha_{1}^{2}}{1-3 \alpha_{1}^{2}} \tag{5.10}
\end{equation*}
$$

provided $3 \alpha_{1}^{2}<1$. These results imply that the kurtosis, $\kappa$, of $y_{t}$ is

$$
\begin{equation*}
\kappa=\frac{E\left(y_{t}^{4}\right)}{\left[E\left(y_{t}^{2}\right)\right]^{2}}=3 \frac{1-\alpha_{1}^{2}}{1-3 \alpha_{1}^{2}}, \tag{5.11}
\end{equation*}
$$

which is always larger than 3 (unless $\alpha_{1}=0$ ), the kurtosis of the normal distribution. Thus, the marginal distribution of the returns, $y_{t}$, is leptokurtic, or has "fat tails." Thus, if $0 \leq \alpha_{1}<1$, the process $y_{t}$ itself is white noise and its unconditional distribution is symmetrically distributed around zero; this distribution is leptokurtic. If, in addition, $3 \alpha_{1}^{2}<1$, the square of the process, $y_{t}^{2}$, follows a causal $\operatorname{AR}(1)$ model with ACF given by $\rho_{y^{2}}(h)=\alpha_{1}^{h} \geq 0$, for all $h>0$. If $3 \alpha_{1} \geq 1$, but $\alpha_{1}<1$, it can be shown that $y_{t}^{2}$ is strictly stationary with infinite variance (see Douc, et al., 2014).

Estimation of the parameters $\alpha_{0}$ and $\alpha_{1}$ of the $\operatorname{ARCH}(1)$ model is typically accomplished by conditional MLE. The conditional likelihood of the data $y_{2}, \ldots ., y_{n}$ given $y_{1}$, is given by

$$
\begin{equation*}
L\left(\alpha_{0}, \alpha_{1} \mid y_{1}\right)=\prod_{t=2}^{n} f_{\alpha_{0}, \alpha_{1}}\left(y_{t} \mid y_{t-1}\right), \tag{5.12}
\end{equation*}
$$

where the density $f_{\alpha_{0}, \alpha_{1}}\left(y_{t} \mid y_{t-1}\right)$ is the normal density specified in (5.5). Hence, the criterion function to be minimized, $l\left(\alpha_{0}, \alpha_{1}\right) \propto-\ln L\left(\alpha_{0}, \alpha_{1} \mid y_{1}\right)$ is given by

$$
\begin{equation*}
l\left(\alpha_{0}, \alpha_{1}\right)=\frac{1}{2} \sum_{t=2}^{n} \ln \left(\alpha_{0}+\alpha_{1} y_{t-1}^{2}\right)+\frac{1}{2} \sum_{t=2}^{n}\left(\frac{y_{t}^{2}}{\alpha_{0}+\alpha_{1} y_{t-1}^{2}}\right) . \tag{5.13}
\end{equation*}
$$

Estimation is accomplished by numerical methods, as described in §3.6.
It is also possible to combine a regression or an ARMA model for the mean with an ARCH model for the errors. For example, a regression with ARCH(1) errors model would have the observations $x_{t}$ as linear function of $p$ regressors, $z_{t}=\left(z_{t 1}, \ldots, z_{t p}\right)^{\prime}$, and $\operatorname{ARCH}(1)$ noise $y_{t}$, say,


Fig. 5.1. ACF and PACF of the squares of the residuals from the AR(1) fit on U.S. GNP.

$$
x_{t}=\boldsymbol{\beta}^{\prime} z_{t}+y_{t}
$$

where $y_{t}$ satisfies (5.3)-(5.4), but, in this case, is unobserved. Similarly, for example, an AR(1) model for data $x_{t}$ exhibiting $\operatorname{ARCH}(1)$ errors would be

$$
x_{t}=\phi_{0}+\phi_{1} x_{t-1}+y_{t} .
$$

These types of models were explored by Weiss (1984).

## Example 5.1 Analysis of U.S. GNP

In Example 3.30, we fit an MA(2) model and an AR(1) model to the U.S. GNP series and we concluded that the residuals from both fits appeared to behave like a white noise process. In Example 3.34 we concluded that the $\operatorname{AR}(1)$ is probably the better model in this case. It has been suggested that the U.S. GNP series has ARCH errors, and in this example, we will investigate this claim. If the GNP noise term is ARCH, the squares of the residuals from the fit should behave like a non-Gaussian $\operatorname{AR}(1)$ process, as pointed out in (5.6). Figure 5.1 shows the ACF and PACF of the squared residuals it appears that there may be some dependence, albeit small, left in the residuals. The figure was generated in R as follows.
gnpgr $=\operatorname{diff}(\log (g n p))$ \# get the growth rate sarima (gnpgr, 1, 0, 0) \# fit an AR(1) acf2(innov^2, 24) \# get ( $p$ )acf of the squared residuals

We used the R package $f$ Garch to fit an $\operatorname{AR}(1)-\operatorname{ARCH}(1)$ model to the U.S. GNP returns with the following results. A partial output is shown; we note that $\operatorname{garch}(1,0)$ specifies an $\mathrm{ARCH}(1)$ in the code below (details later).

| library (fGarch) |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
|  | Estimate Std. Er | Error t | value $\operatorname{Pr}(>t)$ |  |
| mu $\quad 5.278 \mathrm{e}-03$ | $5.278 \mathrm{e}-038.996$ | 6e-04 | 5.8674 .44 | 4e-09 |
| ar1 3.666e-01 | $3.666 \mathrm{e}-01 \quad 7.514$ | $4 \mathrm{e}-02$ | 4.8781 .07 | 7e-06 |
| omega 7.331e-05 | $7.331 \mathrm{e}-059.011$ | 1e-06 | 8.1354 .44 | 4e-16 |
| alpha1 1.945e-01 | $1.945 \mathrm{e}-01$ 9.554 | $4 \mathrm{e}-02$ | 2.0350. | . 0418 |
| Standardised Residuals |  | Tests | Statistic | -Value |
| Jarque-Bera Test | -Bera Test R | Chi^2 | 9.118036 | 0.01047234 |
| Shapiro-Wilk Test | ro-Wilk Test R | W | 0.9842407 | 0.01433690 |
| Ljung-Box Test | Box Test R | Q(10) | 9.874326 | 0.4515875 |
| Ljung-Box Test | Box Test $\quad \mathrm{R}$ | Q(15) | 17.55855 | 0.2865844 |
| Ljung-Box Test | Box Test R | Q(20) | 23.41363 | 0.2689437 |
| Ljung-Box Test | -Box Test $\mathrm{R}^{\wedge} 2$ | Q(10) | 19.2821 | 0.03682246 |
| Ljung-Box Test | -Box Test $\mathrm{R}^{\wedge} 2$ | Q(15) | 33.23648 | 0.004352736 |
| Ljung-Box Test | -Box Test $\mathrm{R}^{\wedge} 2$ | Q(20) | 37.74259 | 0.009518992 |
| LM Arch Test | Test $\quad$ R | TR^2 | 25.41625 | 0.01296901 |

In this example, we obtain $\widehat{\phi}_{0}=.005$ (called mu in the output) and $\widehat{\phi}_{1}=.367$ (called ar1) for the $\operatorname{AR}(1)$ parameter estimates; in Example 3.30 the values were .005 and .347 , respectively. The $\operatorname{ARCH}(1)$ parameter estimates are $\widehat{\alpha}_{0}=0$ (called omega) for the constant and $\widehat{\alpha}_{1}=.195$, which is significant with a p-value of about .04 . There are a number of tests that are performed on the residuals $[\mathrm{R}]$ or the squared residuals [ $\mathrm{R}^{\wedge} 2$ ]. For example, the Jarque-Bera statistic tests the residuals of the fit for normality based on the observed skewness and kurtosis, and it appears that the residuals have some non-normal skewness and kurtosis. The Shapiro-Wilk statistic tests the residuals of the fit for normality based on the empirical order statistics. The other tests, primarily based on the Q-statistic, are used on the residuals and their squares.

The $\operatorname{ARCH}(1)$ model can be extended to the general $\operatorname{ARCH}(m)$ model in an obvious way. That is, (5.3), $y_{t}=\sigma_{t} \epsilon_{t}$, is retained, but (5.4) is extended to

$$
\begin{equation*}
\sigma_{t}^{2}=\alpha_{0}+\alpha_{1} y_{t-1}^{2}+\cdots+\alpha_{m} y_{t-m}^{2} . \tag{5.14}
\end{equation*}
$$

Estimation for $\operatorname{ARCH}(m)$ also follows in an obvious way from the discussion of estimation for $\operatorname{ARCH}(1)$ models. That is, the conditional likelihood of the data $y_{m+1}, \ldots, y_{n}$ given $y_{1}, \ldots, y_{m}$, is given by

$$
\begin{equation*}
L\left(\boldsymbol{\alpha} \mid y_{1}, \ldots, y_{m}\right)=\prod_{t=m+1}^{n} f_{\boldsymbol{\alpha}}\left(y_{t} \mid y_{t-1}, \ldots, y_{t-m}\right), \tag{5.15}
\end{equation*}
$$

where $\boldsymbol{\alpha}=\left(\alpha_{0}, \alpha_{1}, \ldots, \alpha_{m}\right)$ and the conditional densities $f_{\boldsymbol{\alpha}}(\cdot \mid \cdot)$ in (5.15) are normal densities; that is, for $t>m$,

$$
y_{t} \mid y_{t-1}, \ldots, y_{t-m} \sim \mathrm{~N}\left(0, \alpha_{0}+\alpha_{1} y_{t-1}^{2}+\cdots+\alpha_{m} y_{t-m}^{2}\right) .
$$

Another extension of ARCH is the generalized ARCH or GARCH model developed by Bollerslev (1986). For example, a $\operatorname{GARCH}(1,1)$ model retains (5.3), $y_{t}=\sigma_{t} \epsilon_{t}$, but extends (5.4) as follows:

$$
\begin{equation*}
\sigma_{t}^{2}=\alpha_{0}+\alpha_{1} y_{t-1}^{2}+\beta_{1} \sigma_{t-1}^{2} \tag{5.16}
\end{equation*}
$$

Under the condition that $\alpha_{1}+\beta_{1}<1$, using similar manipulations as in (5.6), the GARCH $(1,1)$ model, (5.3) and (5.16), admits a non-Gaussian ARMA $(1,1)$ model for the squared process

$$
\begin{equation*}
y_{t}^{2}=\alpha_{0}+\left(\alpha_{1}+\beta_{1}\right) y_{t-1}^{2}+v_{t}-\beta_{1} v_{t-1} \tag{5.17}
\end{equation*}
$$

where $v_{t}$ is as defined in (5.6). Representation (5.17) follows by writing (5.3) as

$$
\begin{aligned}
y_{t}^{2}-\sigma_{t}^{2} & =\sigma_{t}^{2}\left(\epsilon_{t}^{2}-1\right) \\
\beta_{1}\left(y_{t-1}^{2}-\sigma_{t-1}^{2}\right) & =\beta_{1} \sigma_{t-1}^{2}\left(\epsilon_{t-1}^{2}-1\right),
\end{aligned}
$$

subtracting the second equation from the first, and using the fact that, from (5.16), $\sigma_{t}^{2}-\beta_{1} \sigma_{t-1}^{2}=\alpha_{0}+\alpha_{1} y_{t-1}^{2}$, on the left-hand side of the result. The $\operatorname{GARCH}(m, r)$ model retains (5.3) and extends (5.16) to

$$
\begin{equation*}
\sigma_{t}^{2}=\alpha_{0}+\sum_{j=1}^{m} \alpha_{j} y_{t-j}^{2}+\sum_{j=1}^{r} \beta_{j} \sigma_{t-j}^{2} \tag{5.18}
\end{equation*}
$$

Conditional maximum likelihood estimation of the $\operatorname{GARCH}(m, r)$ model parameters is similar to the $\operatorname{ARCH}(m)$ case, wherein the conditional likelihood, (5.15), is the product of $\mathrm{N}\left(0, \sigma_{t}^{2}\right)$ densities with $\sigma_{t}^{2}$ given by (5.18) and where the conditioning is on the first $\max (m, r)$ observations, with $\sigma_{1}^{2}=\cdots=\sigma_{r}^{2}=0$. Once the parameter estimates are obtained, the model can be used to obtain one-step-ahead forecasts of the volatility, say $\widehat{\sigma}_{t+1}^{2}$, given by

$$
\begin{equation*}
\widehat{\sigma}_{t+1}^{2}=\widehat{\alpha}_{0}+\sum_{j=1}^{m} \widehat{\alpha}_{j} y_{t+1-j}^{2}+\sum_{j=1}^{r} \widehat{\beta}_{j} \widehat{\sigma}_{t+1-j}^{2} \tag{5.19}
\end{equation*}
$$

We explore these concepts in the following example.

## Example 5.2 GARCH Analysis of the NYSE Returns

As previously mentioned, the daily returns of the NYSE shown in Figure 1.4 exhibit classic GARCH features. We used the R fGarch package to fit a $\operatorname{GARCH}(1,1)$ model to the series with the following results:

```
library(fGarch)
nyse = astsa::nyse # fGarch has one too
summary(nyse.g <- garchFit(~garch(1,1), nyse))
            Estimate Std. Error t value Pr(>t)
    mu 7.369e-04 1.786e-04 4.126 3.69e-05
    omega 6.542e-06 1.455e-06 4.495 6.94e-06
    alpha1 1.141e-01 1.604e-02 7.114 1.13e-12
    beta1 8.061e-01 2.973e-02 27.112 < 2e-16
Standardised Residuals Tests:
    Statistic p-Value
    Jarque-Bera Test R Chi^2 3628.415 Q
    Shapiro-Wilk Test R W 0.9515562 Q
```



Fig. 5.2. GARCH predictions of the NYSE volatility, $\pm 2 \widehat{\sigma}_{t}$, displayed as dashed lines.

| Ljung-Box Test | $R$ | $Q(10)$ | 29.69242 | 0.0009616813 |
| :--- | :--- | :--- | :--- | :--- |
| Ljung-Box Test | $R$ | $Q(15)$ | 30.50938 | 0.01021164 |
| Ljung-Box Test | $R$ | $Q(20)$ | 32.81143 | 0.03538324 |
| Ljung-Box Test | $R^{\wedge} 2$ | $Q(10)$ | 3.510505 | 0.9667405 |
| Ljung-Box Test | $R^{\wedge} 2$ | $Q(15)$ | 4.408852 | 0.9960585 |
| Ljung-Box Test | $R^{\wedge} 2$ | $Q(20)$ | 6.68935 | 0.9975864 |
| LM Arch Test | $R$ | $T R^{\wedge} 2$ | 3.967784 | 0.9840107 |

To explore the GARCH predictions of volatility, we calculated and plotted the 100 observations from the middle of the data (which includes the October 19, 1987 crash) along with the one-step-ahead predictions of the corresponding volatility, $\sigma_{t}^{2}$. The results are displayed as the data $\pm 2 \widehat{\sigma}_{t}$ as a dashed line surrounding the data in Figure 5.2.

```
u = nyse.g@sigma.t
plot(window(nyse, start=900, end=1000), ylim=c(-.22,.2), ylab="NYSE Returns")
lines(window(nyse-2*u, start=900, end=1000), lty=2, col=4)
lines(window(nyse+2*u, start=900, end=1000), lty=2, col=4)
```

Some key points can be gleaned from the examples of this section. First, it is apparent that the conditional distribution of the returns is rarely normal. fGarch allows for various distributions to be fit to the data; see the help file for information. Some drawbacks of the GARCH model are: (i) the model assumes positive and negative returns have the same effect because volatility depends on squared returns; (ii) the model is restrictive because of the tight constraints on the model parameters (e.g., for an $\operatorname{ARCH}(1), 0 \leq \alpha_{1}^{2}<\frac{1}{3}$ ); (iii) the likelihood is flat unless $n$ is very large; (iv) the model tends to overpredict volatility because it responds slowly to large isolated returns.

Various extensions to the original model have been proposed to overcome some of the shortcomings we have just mentioned. For example, we have already discussed the fact that the S-PLUS Garch module will fit some non-normal, albeit symmetric, distributions. For asymmetric return dynamics, one can use the EGARCH (exponential GARCH) model, which is a complex model that has different components for
positive returns and for negative returns. In the case of persistence in volatility, the integrated GARCH (IGARCH) model may be used. Recall (5.17) where we showed the $\operatorname{GARCH}(1,1)$ model can be written as

$$
y_{t}^{2}=\alpha_{0}+\left(\alpha_{1}+\beta_{1}\right) y_{t-1}^{2}+v_{t}-\beta_{1} v_{t-1}
$$

and $y_{t}^{2}$ is stationary if $\alpha_{1}+\beta_{1}<1$. The IGARCH model sets $\alpha_{1}+\beta_{1}=1$, in which case the $\operatorname{IGARCH}(1,1)$ model is

$$
y_{t}=\sigma_{t} \epsilon_{t} \quad \text { and } \quad \sigma_{t}^{2}=\alpha_{0}+\left(1-\beta_{1}\right) y_{t-1}^{2}+\beta_{1} \sigma_{t-1}^{2}
$$

There are many different extensions to the basic ARCH model that were developed to handle the various situations noticed in practice. Interested readers might find the general discussions in Engle et al. (1994) and Shephard (1996) worthwhile reading. Also, Gouriéroux (1997) gives a detailed presentation of ARCH and related models with financial applications and contains an extensive bibliography. Two excellent texts on financial time series analysis are Chan (2002) and Tsay (2002).

### 5.2 Long Memory ARMA and Fractional Differencing

The conventional ARMA $(p, q)$ process is often referred to as a short-memory process because the coefficients in the representation

$$
x_{t}=\sum_{j=0}^{\infty} \psi_{j} w_{t-j}
$$

obtained by solving

$$
\phi(z) \psi(z)=\theta(z),
$$

are dominated by exponential decay. As pointed out in Chapter 3, this result implies the ACF of the short memory process $\rho(h) \rightarrow 0$ exponentially fast as $h \rightarrow \infty$. When the sample ACF of a time series decays slowly, the advice given in Chapter 3 has been to difference the series until it seems stationary. Following this advice with the glacial varve series first presented in Example 3.20 leads to the first difference of the logarithms of the data being represented as a first-order moving average. In Example 3.32, further analysis of the residuals leads to fitting an $\operatorname{ARIMA}(1,1,1)$ model,

$$
\nabla x_{t}=\phi \nabla x_{t-1}+w_{t}+\theta w_{t-1}
$$

where we understand $x_{t}$ is the log-transformed varve series. In particular, the estimates of the parameters (and the standard errors) were $\widehat{\phi}=.23(.05), \widehat{\theta}=-.89(.03)$, and $\widehat{\sigma}_{w}^{2}=.23$. The use of the first difference $\nabla x_{t}=(1-B) x_{t}$ can be too severe a modification in the sense that the nonstationary model might represent an overdifferencing of the original process.

Long memory (or persistent) time series were considered in Hosking (1981) and Granger and Joyeux (1980) as intermediate compromises between the short memory

## log(varve)



Fig. 5.3. Sample ACF of the log transformed varve series.

ARMA type models and the fully integrated nonstationary processes in the BoxJenkins class. The easiest way to generate a long memory series is to think of using the difference operator $(1-B)^{d}$ for fractional values of $d$, say, $0<d<.5$, so a basic long memory series gets generated as

$$
\begin{equation*}
(1-B)^{d} x_{t}=w_{t}, \tag{5.20}
\end{equation*}
$$

where $w_{t}$ still denotes white noise with variance $\sigma_{w}^{2}$. The fractionally differenced series (5.20), for $|d|<.5$, is often called fractional noise (except when $d$ is zero). Now, $d$ becomes a parameter to be estimated along with $\sigma_{w}^{2}$. Differencing the original process, as in the Box-Jenkins approach, may be thought of as simply assigning a value of $d=1$. This idea has been extended to the class of fractionally integrated ARMA, or ARFIMA models, where $-.5<d<.5$; when $d$ is negative, the term antipersistent is used. Long memory processes occur in hydrology (see Hurst, 1951, and McLeod and Hipel, 1978) and in environmental series, such as the varve data we have previously analyzed, to mention a few examples. Long memory time series data tend to exhibit sample autocorrelations that are not necessarily large (as in the case of $d=1$ ), but persist for a long time. Figure 5.3 shows the sample ACF, to lag 100, of the log-transformed varve series, which exhibits classic long memory behavior:
$\mathrm{u}=\mathrm{acf}(\log$ (varve), 100 , plot=FALSE)
plot(u[1:100], ylim=c (-.1,1), main="log(varve)") \# get rid of lag 0
To investigate its properties, we can use the binomial expansion $(d>-1)$ to write

$$
\begin{equation*}
w_{t}=(1-B)^{d} x_{t}=\sum_{j=0}^{\infty} \pi_{j} B^{j} x_{t}=\sum_{j=0}^{\infty} \pi_{j} x_{t-j} \tag{5.21}
\end{equation*}
$$

where

$$
\begin{equation*}
\pi_{j}=\frac{\Gamma(j-d)}{\Gamma(j+1) \Gamma(-d)} \tag{5.22}
\end{equation*}
$$

with $\Gamma(x+1)=x \Gamma(x)$ being the gamma function. Similarly $(d<1)$, we can write

$$
\begin{equation*}
x_{t}=(1-B)^{-d} w_{t}=\sum_{j=0}^{\infty} \psi_{j} B^{j} w_{t}=\sum_{j=0}^{\infty} \psi_{j} w_{t-j} \tag{5.23}
\end{equation*}
$$

where

$$
\begin{equation*}
\psi_{j}=\frac{\Gamma(j+d)}{\Gamma(j+1) \Gamma(d)} \tag{5.24}
\end{equation*}
$$

When $|d|<.5$, the processes (5.21) and (5.23) are well-defined stationary processes (see Brockwell and Davis, 1991, for details). In the case of fractional differencing, however, the coefficients satisfy $\sum \pi_{j}^{2}<\infty$ and $\sum \psi_{j}^{2}<\infty$ as opposed to the absolute summability of the coefficients in ARMA processes.

Using the representation (5.23)-(5.24), and after some nontrivial manipulations, it can be shown that the ACF of $x_{t}$ is

$$
\begin{equation*}
\rho(h)=\frac{\Gamma(h+d) \Gamma(1-d)}{\Gamma(h-d+1) \Gamma(d)} \sim h^{2 d-1} \tag{5.25}
\end{equation*}
$$

for large $h$. From this we see that for $0<d<.5$

$$
\sum_{h=-\infty}^{\infty}|\rho(h)|=\infty
$$

and hence the term long memory.
In order to examine a series such as the varve series for a possible long memory pattern, it is convenient to look at ways of estimating $d$. Using (5.22) it is easy to derive the recursions

$$
\begin{equation*}
\pi_{j+1}(d)=\frac{(j-d) \pi_{j}(d)}{(j+1)} \tag{5.26}
\end{equation*}
$$

for $j=0,1, \ldots$, with $\pi_{0}(d)=1$. Maximizing the joint likelihood of the errors under normality, say, $w_{t}(d)$, will involve minimizing the sum of squared errors

$$
Q(d)=\sum w_{t}^{2}(d)
$$

The usual Gauss-Newton method, described in §3.6, leads to the expansion

$$
w_{t}(d)=w_{t}\left(d_{0}\right)+w_{t}^{\prime}\left(d_{0}\right)\left(d-d_{0}\right)
$$

where

$$
w_{t}^{\prime}\left(d_{0}\right)=\left.\frac{\partial w_{t}}{\partial d}\right|_{d=d_{0}}
$$

and $d_{0}$ is an initial estimate (guess) at to the value of $d$. Setting up the usual regression leads to


Fig. 5.4. Coefficients $\pi_{j}(.384), j=1,2, \ldots, 30$ in the representation (5.26).

$$
\begin{equation*}
d=d_{0}-\frac{\sum_{t} w_{t}^{\prime}\left(d_{0}\right) w_{t}\left(d_{0}\right)}{\sum_{t} w_{t}^{\prime}\left(d_{0}\right)^{2}} . \tag{5.27}
\end{equation*}
$$

The derivatives are computed recursively by differentiating (5.26) successively with respect to $d$ : $\pi_{j+1}^{\prime}(d)=\left[(j-d) \pi_{j}^{\prime}(d)-\pi_{j}(d)\right] /(j+1)$, where $\pi_{0}^{\prime}(d)=0$. The errors are computed from an approximation to (5.21), namely,

$$
\begin{equation*}
w_{t}(d)=\sum_{j=0}^{t} \pi_{j}(d) x_{t-j} . \tag{5.28}
\end{equation*}
$$

It is advisable to omit a number of initial terms from the computation and start the sum, (5.27), at some fairly large value of $t$ to have a reasonable approximation.

## Example 5.3 Long Memory Fitting of the Glacial Varve Series

We consider analyzing the glacial varve series discussed in Example 2.6 and Example 3.20. Figure 2.6 shows the original and log-transformed series (which we denote by $x_{t}$ ). In Example 3.32, we noted that $x_{t}$ could be modeled as an $\operatorname{ARIMA}(1,1,1)$ process. We fit the fractionally differenced model, (5.20), to the mean-adjusted series, $x_{t}-\bar{x}$. Applying the Gauss-Newton iterative procedure previously described, starting with $d=.1$ and omitting the first 30 points from the computation, leads to a final value of $d=.384$, which implies the set of coefficients $\pi_{j}(.384)$, as given in Figure 5.4 with $\pi_{0}(.384)=1$. We can compare roughly the performance of the fractional difference operator with the ARIMA model by examining the autocorrelation functions of the two residual series as shown in Figure 5.5. The ACFs of the two residual series are roughly comparable with the white noise model.

To perform this analysis in R, first download and install the fracdiff package. Then use


Fig. 5.5. ACF of residuals from the $\operatorname{ARIMA}(1,1,1)$ fit to the logged varve series (top) and of the residuals from the long memory model fit, $(1-B)^{d} x_{t}=w_{t}$, with $d=.384$ (bottom).

```
p = rep(1,31)
for (k in 1:30){ p[k+1] = (k-varve.fd$d)*p[k]/(k+1) }
plot(1:30, p[-1], ylab=expression(pi(d)), xlab="Index", type="h")
res.fd = diffseries(log(varve), varve.fd$d) # frac diff resids
res.arima = resid(arima(log(varve), order=c(1,1,1))) # arima resids
par(mfrow=c(2,1))
acf(res.arima, 100, xlim=c(4,97), ylim=c(-.2,.2), main="")
acf(res.fd, 100, xlim=c(4,97), ylim=c(-.2,.2), main="")
```

The R package uses a truncated maximum likelihood procedure that was discussed in Haslett and Raftery (1989), which is a little more elaborate than simply zeroing out initial values. The default truncation value in R is $M=100$. In the default case, the estimate is $\widehat{d}=.37$ with approximately the same (questionable) standard error.

Forecasting long memory processes is similar to forecasting ARIMA models. That is, (5.21) and (5.26) can be used to obtain the truncated forecasts

$$
\begin{equation*}
\widetilde{x}_{n+m}^{n}=-\sum_{j=1}^{n} \pi_{j}(\widehat{d}) \widetilde{x}_{n+m-j}^{n}, \tag{5.29}
\end{equation*}
$$

for $m=1,2, \ldots$. Error bounds can be approximated by using

$$
\begin{equation*}
P_{n+m}^{n}=\widehat{\sigma}_{w}^{2}\left(\sum_{j=0}^{m-1} \psi_{j}^{2}(\widehat{d})\right) \tag{5.30}
\end{equation*}
$$

where, as in (5.26),

$$
\begin{equation*}
\psi_{j}(\widehat{d})=\frac{(j+\widehat{d}) \psi_{j}(\widehat{d})}{(j+1)} \tag{5.31}
\end{equation*}
$$

with $\psi_{0}(\widehat{d})=1$.
No obvious short memory ARMA-type component can be seen in the ACF of the residuals from the fractionally differenced varve series shown in Figure 5.5. It is natural, however, that cases will exist in which substantial short memory-type components will also be present in data that exhibits long memory. Hence, it is natural to define the general $\operatorname{ARFIMA}(p, d, q),-.5<d<.5$ process as

$$
\begin{equation*}
\phi(B) \nabla^{d}\left(x_{t}-\mu\right)=\theta(B) w_{t}, \tag{5.32}
\end{equation*}
$$

where $\phi(B)$ and $\theta(B)$ are as given in Chapter 3. Writing the model in the form

$$
\begin{equation*}
\phi(B) \pi_{d}(B)\left(x_{t}-\mu\right)=\theta(B) w_{t} \tag{5.33}
\end{equation*}
$$

makes it clear how we go about estimating the parameters for the more general model. Forecasting for the $\operatorname{ARFIMA}(p, d, q)$ series can be easily done, noting that we may equate coefficients in

$$
\begin{equation*}
\phi(z) \psi(z)=(1-z)^{-d} \theta(z) \tag{5.34}
\end{equation*}
$$

and

$$
\begin{equation*}
\theta(z) \pi(z)=(1-z)^{d} \phi(z) \tag{5.35}
\end{equation*}
$$

to obtain the representations

$$
x_{t}=\mu+\sum_{j=0}^{\infty} \psi_{j} w_{t-j}
$$

and

$$
w_{t}=\sum_{j=0}^{\infty} \pi_{j}\left(x_{t-j}-\mu\right) .
$$

We then can proceed as discussed in (5.29) and (5.30).
Comprehensive treatments of long memory time series models are given in the texts by Beran (1994), Palma (2007), and Robinson (2003), and it should be noted that several other techniques for estimating the parameters, especially, the long memory parameter, can be developed in the frequency domain. In this case, we may think of the equations as generated by an infinite order autoregressive series with coefficients $\pi_{j}$ given by (5.26). Using the same approach as before, we obtain

$$
\begin{align*}
f_{x}(\omega) & =\frac{\sigma_{w}^{2}}{\left|\sum_{k=0}^{\infty} \pi_{k} \epsilon^{-2 \pi i k \omega}\right|^{2}}  \tag{5.36}\\
& =\sigma_{w}^{2}\left|1-e^{-2 \pi i \omega}\right|^{-2 d}=\left[4 \sin ^{2}(\pi \omega)\right]^{-d} \sigma_{w}^{2}
\end{align*}
$$

as equivalent representations of the spectrum of a long memory process. The long memory spectrum approaches infinity as the frequency $\omega \rightarrow 0$.

The main reason for defining the Whittle approximation to the log likelihood is to propose its use for estimating the parameter $d$ in the long memory case as an alternative to the time domain method previously mentioned. The time domain approach is useful because of its simplicity and easily computed standard errors. One may also use an exact likelihood approach by developing an innovations form of the likelihood as in Brockwell and Davis (1991).

For the approximate approach using the Whittle likelihood, we consider using the approach of Fox and Taqqu (1986) who showed that maximizing the Whittle log likelihood leads to a consistent estimator with the usual asymptotic normal distribution that would be obtained by treating the Whittle likelihood as a conventional log likelihood (see also Dahlhaus, 1989; Robinson, 1995; Hurvich et al., 1998). Unfortunately, the periodogram ordinates are not asymptotically independent (Hurvich and Beltrao, 1993), although a quasi-likelihood in the form of the Whittle approximation works well and has good asymptotic properties.

To see how this would work for the purely long memory case, write the long memory spectrum as

$$
\begin{equation*}
f_{x}\left(\omega_{k} ; d, \sigma_{w}^{2}\right)=\sigma_{w}^{2} g_{k}^{-d} \tag{5.37}
\end{equation*}
$$

where

$$
\begin{equation*}
g_{k}=4 \sin ^{2}\left(\pi \omega_{k}\right) \tag{5.38}
\end{equation*}
$$

Then, differentiating the log likelihood, say,

$$
\begin{equation*}
\ln L\left(\boldsymbol{x} ; d, \sigma_{w}^{2}\right) \approx-m \ln \sigma_{w}^{2}+d \sum_{k=1}^{m} \ln g_{k}-\frac{1}{\sigma_{w}^{2}} \sum_{k=1}^{m} g_{k}^{d} I\left(\omega_{k}\right) \tag{5.39}
\end{equation*}
$$

at $m=n / 2-1$ frequencies and solving for $\sigma_{w}^{2}$ yields

$$
\begin{equation*}
\sigma_{w}^{2}(d)=\frac{1}{m} \sum_{k=1}^{m} g_{k}^{d} I\left(\omega_{k}\right) \tag{5.40}
\end{equation*}
$$

as the approximate maximum likelihood estimator for the variance parameter. To estimate $d$, we can use a grid search of the concentrated log likelihood

$$
\begin{equation*}
\ln L(\boldsymbol{x} ; d) \approx-m \ln \sigma_{w}^{2}(d)+d \sum_{k=1}^{m} \ln g_{k}-m \tag{5.41}
\end{equation*}
$$

over the interval ( $-.5, .5$ ), followed by a Newton-Raphson procedure to convergence.

## Example 5.4 Long Memory Spectra for the Varve Series

In Example 5.3, we fit a long memory model to the glacial varve data via time domain methods. Fitting the same model using frequency domain methods and the Whittle approximation above gives $\widehat{d}=.380$, with an estimated standard error of .028. The earlier time domain method gave $\widehat{d}=.384$ with $M=30$ and $\widehat{d}=.370$ with $M=100$. Both estimates obtained via time domain methods had a standard error of about $4.6 \times 10^{-6}$, which seems implausible. The error variance estimate in this case is $\widehat{\sigma}_{w}^{2}=.2293$; in Example 5.3, we could have used var(res.fd) as an estimate, in which case we obtain .2298 . The R code to perform this analysis is


Fig. 5.6. Long Memory $(d=.380)$ [solid line] and autoregressive $\operatorname{AR}(8)$ [dashed line] spectral estimators for the paleoclimatic glacial varve series.

```
series = log(varve) # specify series to be analyzed
d0 = .1 # initial value of d
n.per = nextn(length(series))
m = (n.per)/2 - 1
per = abs(fft(series-mean(series))[-1])^2 # remove 0 freq
per = per/n.per # and scale the peridogram
g = 4*(sin(pi*((1:m)/n.per))^2)
# Function to calculate -log.likelihood
whit.like = function(d){
    g.d=g^d
    sig2 = (sum(g.d*per[1:m])/m)
    log.like = m*log(sig2) - d*sum(log(g)) + m
    return(log.like) }
# Estimation (?optim for details - output not shown)
(est = optim(dQ, whit.like, gr=NULL, method="L-BFGS-B", hessian=TRUE,
    lower=-.5, upper=.5, control=list(trace=1,REPORT=1)))
# Results: d.hat = .380, se(dhat) = .028, and sig2hat = . }22
cat("d.hat =", est$par, "se(dhat) = ",1/sqrt(est$hessian),"\n")
g.dhat = g^est$par; sig2 = sum(g.dhat*per[1:m])/m
cat("sig2hat =",sig2,"\n")
```

One might also consider fitting an autoregressive model to these data using a procedure similar to that used in Example 4.12. Following this approach gave an autoregressive model with $p=8$ and $\widehat{\boldsymbol{\phi}}=(.34, .11, .04, .09, .08, .08, .02, .09)^{\prime}$, with $\widehat{\sigma}_{w}^{2}=.2267$ as the error variance. The two log spectra are plotted in Figure 5.6 for $\omega>0$, and we note that long memory spectrum will eventually become infinite, whereas the $\operatorname{AR}(8)$ spectrum is finite at $\omega=0$. The R code used for this part of the example (assuming the previous values have been retained) is

```
u = spec.ar(log(varve), plot=FALSE) # produces AR(8)
g = 4*(sin(pi*((1:500)/2000))^2)
```

```
fhat = sig2*g^{-est$par} # long memory spectral estimate
plot(1:500/2000, log(fhat), type="l", ylab="log(spectrum)", xlab="frequency")
lines(u$freq[1:250], log(u$spec[1:250]), lty="dashed")
ar.mle(log(varve)) # to get AR(8) estimates
```

Often, time series are not purely long memory. A common situation has the long memory component multiplied by a short memory component, leading to an alternate version of (5.37) of the form

$$
\begin{equation*}
f_{x}\left(\omega_{k} ; d, \theta\right)=g_{k}^{-d} f_{0}\left(\omega_{k} ; \boldsymbol{\theta}\right), \tag{5.42}
\end{equation*}
$$

where $f_{0}\left(\omega_{k} ; \theta\right)$ might be the spectrum of an autoregressive moving average process with vector parameter $\boldsymbol{\theta}$, or it might be unspecified. If the spectrum has a parametric form, the Whittle likelihood can be used. However, there is a substantial amount of semiparametric literature that develops the estimators when the underlying spectrum $f_{0}(\omega ; \boldsymbol{\theta})$ is unknown. A class of Gaussian semi-parametric estimators simply uses the same Whittle likelihood (5.41), evaluated over a sub-band of low frequencies, say $m^{\prime}=\sqrt{n}$. There is some latitude in selecting a band that is relatively free from low frequency interference due to the short memory component in (5.42).

Geweke and Porter-Hudak (1983) developed an approximate method for estimating $d$ based on a regression model, derived from (5.41). Note that we may write a simple equation for the logarithm of the spectrum as

$$
\begin{equation*}
\ln f_{x}\left(\omega_{k} ; d\right)=\ln f_{0}\left(\omega_{k} ; \boldsymbol{\theta}\right)-d \ln \left[4 \sin ^{2}\left(\pi \omega_{k}\right)\right], \tag{5.43}
\end{equation*}
$$

with the frequencies $\omega_{k}=k / n$ restricted to a range $k=1,2, \ldots, m^{\prime}$ near the zero frequency with $m^{\prime}=\sqrt{n}$ as the recommended value. Relationship (5.43) suggests using a simple linear regression model of the form,

$$
\begin{equation*}
\ln I\left(\omega_{k}\right)=\beta_{0}-d \ln \left[4 \sin ^{2}\left(\pi \omega_{k}\right)\right]+e_{k} \tag{5.44}
\end{equation*}
$$

for the periodogram to estimate the parameters $\sigma_{w}^{2}$ and $d$. In this case, one performs least squares using $\ln I\left(\omega_{k}\right)$ as the dependent variable, and $\ln \left[4 \sin ^{2}\left(\pi \omega_{k}\right)\right]$ as the independent variable for $k=1,2, \ldots, m$. The resulting slope estimate is then used as an estimate of $-d$. For a good discussion of various alternative methods for selecting $m$, see Hurvich and Deo (1999). The R package fracdiff also provides this method via the command fdGPH() ; see the help file for further information.

One of the above two procedures works well for estimating the long memory component but there will be cases (such as ARFIMA) where there will be a parameterized short memory component $f_{0}\left(\omega_{k} ; \boldsymbol{\theta}\right)$ that needs to be estimated. If the spectrum is highly parameterized, one might estimate using the Whittle log likelihood (5.38) and

$$
f_{x}\left(\omega_{k} ; \boldsymbol{\theta}\right)=g_{k}^{-d} f_{0}\left(\omega_{k} ; \boldsymbol{\theta}\right)
$$

and jointly estimating the parameters $d$ and $\boldsymbol{\theta}$ using the Newton-Raphson method. If we are interested in a nonparametric estimator, using the conventional smoothed spectral estimator for the periodogram, adjusted for the long memory component, say $g_{k}^{d} I\left(\omega_{k}\right)$ might be a possible approach.

### 5.3 Unit Root Testing

As discussed in the previous section, the use of the first difference $\nabla x_{t}=(1-B) x_{t}$ can be too severe a modification in the sense that the nonstationary model might represent an overdifferencing of the original process.

Consider a causal AR(1) process (we assume throughout this section that the noise is Gaussian),

$$
\begin{equation*}
x_{t}=\phi x_{t-1}+w_{t} . \tag{5.45}
\end{equation*}
$$

A unit root test provides a way to test whether (5.45) is a random walk (the null case) as opposed to a causal process (the alternative). That is, it provides a procedure for testing

$$
H_{0}: \phi=1 \quad \text { versus } \quad H_{1}:|\phi|<1 .
$$

To see if it is reasonable to assume $\phi-1=0$, an obvious test statistic would be to consider ( $\widehat{\phi}-1$ ), appropriately normalized, in the hope to develop an asymptotically normal test statistic, where $\widehat{\phi}$ is one of the optimal estimators discussed in Chapter 3, §3.4. Based on Example 3.21, one might think $\sqrt{n}(\widehat{\phi}-1)$ is approximately normal. Unfortunately, the theory of $\S 3.4$ will not work in the null case because the process is nonstationary. However, the test statistic

$$
U=n(\widehat{\phi}-1)
$$

can be used, and it is known as the unit root or Dickey-Fuller (DF) statistic (see Fuller, 1996), although the actual DF test statistic is normalized a little differently. In this case, the distribution of the test statistic does not have a closed form and quantiles of the distribution must be computed by numerical approximation or by simulation. The R package tseries provides this test along with more general tests that we mention briefly.

Toward a more general model, we note that the DF test was established by noting that if $x_{t}=\phi x_{t-1}+w_{t}$, then $\nabla x_{t}=(\phi-1) x_{t-1}+w_{t}=\gamma x_{t-1}+w_{t}$, and one could test $H_{0}: \gamma=0$ by regressing $\nabla x_{t}$ on $x_{t-1}$. They formed a Wald statistic and derived its limiting distribution. The test was extended to accommodate $\operatorname{AR}(p)$ models, $x_{t}=\sum_{j=1}^{p} \phi_{j} x_{t-j}+w_{t}$, as follows. Subtract $x_{t-1}$ from the model to obtain

$$
\begin{equation*}
\nabla x_{t}=\gamma x_{t-1}+\sum_{j=1}^{p-1} \psi_{j} \nabla x_{t-j}+w_{t} \tag{5.46}
\end{equation*}
$$

where $\gamma=\sum_{j=1}^{p} \phi_{j}-1$ and $\psi_{j}=-\sum_{j=i}^{p} \phi_{i}$ for $j=2, \ldots, p$. For a quick check of (5.46) when $p=2$, note that $x_{t}=\left(\phi_{1}+\phi_{2}\right) x_{t-1}-\phi_{2}\left(x_{t-1}-x_{t-2}\right)+w_{t}$; now subtract $x_{t-1}$ from both sides. To test the hypothesis that the process has a unit root at 1 (i.e., the AR polynoimial $\phi(z)=0$ when $z=1$ ), we can test $H_{0}: \gamma=0$ by estimating $\gamma$ in the regression of $\nabla x_{t}$ on $x_{t-1}, \nabla x_{t-1}, \ldots, \nabla x_{t-p+1}$, and forming a Wald test based on $t_{\gamma}=\widehat{\gamma} / \operatorname{se}(\widehat{\gamma})$. This test leads to the so-called augmented Dickey-Fuller test (ADF). While the calculations for obtaining the asymptotic null distribution change, the basic ideas and machinery remain the same as in the simple case. The choice of $p$ is crucial,
and we will discuss some suggestions in the example. For $\operatorname{ARMA}(p, q)$ models, the ADF test can be used by assuming $p$ is large enough to capture the essential correlation structure; another alternative is the Phillips-Perron (PP) test, which differs from the ADF tests mainly in how they deal with serial correlation and heteroskedasticity in the errors.

One can extend the model to include a constant, or even non-stochastic trend. For example, consider the model

$$
x_{t}=\beta_{0}+\beta_{1} t+\phi x_{t-1}+w_{t} .
$$

If we assume $\beta_{1}=0$, then under the null hypothesis, $\phi=1$, the process is a random walk with drift $\beta_{0}$. Under the alternate hypothesis, the process is a causal $\operatorname{AR}(1)$ with mean $\mu_{x}=\beta_{0}(1-\phi)$. If we cannot assume $\beta_{1}=0$, then the interest here is testing the null that $\left(\beta_{1}, \phi\right)=(0,1)$, simultaneously, versus the alternative that $\beta_{1} \neq 0$ and $|\phi|<1$. In this case, the null hypothesis is that the process is a random walk with drift, versus the alternative hypothesis that the process is stationary around a global trend (consider the global temperature series examined in Example 2.1).

## Example 5.5 Testing Unit Roots in the Glacial Varve Series

In this example we use the R package tseries to test the null hypothesis that the $\log$ of the glacial varve series has a unit root, versus the alternate hypothesis that the process is stationary. We test the null hypothesis using the available DF, ADF and PP tests; note that in each case, the general regression equation incorporates a constant and a linear trend. In the ADF test, the default number of AR components included in the model, say $k$, is $\llbracket(n-1)^{\frac{1}{3}} \rrbracket$, which corresponds to the suggested upper bound on the rate at which the number of lags, $k$, should be made to grow with the sample size for the general $\operatorname{ARMA}(p, q)$ setup. For the PP test, the default value of $k$ is $\llbracket .04 n^{\frac{1}{4}} \rrbracket$.
library(tseries)
adf.test(log(varve), k=0) \# DF test
Dickey-Fuller $=-12.8572$, Lag order $=0, \mathrm{p}$-value $<0.01$ alternative hypothesis: stationary

```
adf.test(log(varve)) # ADF test
    Dickey-Fuller = -3.5166, Lag order = 8, p-value = 0.04071
        alternative hypothesis: stationary
pp.test(log(varve)) # PP test
    Dickey-Fuller Z(alpha) = -304.5376,
        Truncation lag parameter = 6, p-value < 0.01
        alternative hypothesis: stationary
```

In each test, we reject the null hypothesis that the logged varve series has a unit root. The conclusion of these tests supports the conclusion of the previous section that the logged varve series is long memory rather than integrated.

## Problems

5.1 Weekly crude oil spot prices in dollars per barrel are in oil. Investigate whether the growth rate of the weekly oil price exhibits GARCH behavior. If so, fit an appropriate model to the growth rate.
5.2 The stats package of $R$ contains the daily closing prices of four major European stock indices; type help(EuStockMarkets) for details. Fit a GARCH model to the returns of one of these series and discuss your findings. (Note: The data set contains actual values, and not returns. Hence, the data must be transformed prior to the model fitting.)
5.3 The data set arf is 1000 simulated observations from an $\operatorname{ARFIMA}(1,1,0)$ model with $\phi=.75$ and $d=.4$.
(a) Plot the data and comment.
(b) Plot the ACF and PACF of the data and comment.
(c) Estimate the parameters and test for the significance of the estimates $\widehat{\phi}$ and $\widehat{d}$.
(d) Explain why, using the results of parts (a) and (b), it would seem reasonable to difference the data prior to the analysis. That is, if $x_{t}$ represents the data, explain why we might choose to fit an ARMA model to $\nabla x_{t}$.
(e) Plot the ACF and PACF of $\nabla x_{t}$ and comment.
(f) Fit an ARMA model to $\nabla x_{t}$ and comment.
5.4 Compute the sample ACF of the absolute values of the NYSE returns displayed in Figure 1.4 up to lag 200, and comment on whether the ACF indicates long memory. Fit an ARFIMA model to the absolute values and comment.
5.5 Plot the global temperature series, gtemp, and then test whether there is a unit root versus the alternative that the process is stationary using the three tests, DF, ADF, and PP, discussed in Example 5.5. Comment.
5.6 Plot the GNP series, gnp, and then test for a unit root against the alternative that the process is explosive. State your conclusion.
5.7 Verify (5.46).

## Appendix $\mathbf{R}$

## R Supplement

## R. 1 First Things First

The website for the text is http://www.stat.pitt.edu/stoffer/tsa3/. If you do not already have R, point your browser to the Comprehensive R Archive Network (CRAN), http://cran.r-project.org/ and download and install it. The installation includes help files and some user manuals. You can find helpful tutorials by following CRAN's link to Contributed Documentation. If you are new to R/S-PLUS, then $R$ for Beginners by Emmanuel Paradis is a great introduction. There is also a lot of advice out there in cyberspace, but some of it will be outdated because R goes through many revisions. Also, beginners will find RStudio (https://www.rstudio.com/) helpful for using R. [Advice: If you've never used a command based program, get $R$ studio.]

## R. 2 ASTSA (astsa)

There is an R package for the text called astsa (Applied Statistical Time Series Analysis), which was the name of the software distributed with the first and second editions of this text, and the original version, Shumway (1988). The package can be obtained from CRAN and its mirrors in the usual way. To download and install astsa, start R and type
install.packages("astsa")
You will be asked to choose the closest CRAN mirror to you. As with all packages, you have to load astsa before you use it by issuing the command
require(astsa)
All the data are loaded when the package is loaded. If you create a .First function as follows,

```
.First <- function(){require(astsa)}
```

and save the workspace when you quit, astsa will be loaded at every start until you change .First.

R is not consistent with help files across different operating systems. The best help system is the html help, which can be started issuing the command help.start() and then following the Packages link to astsa. In Windows, you can type help(package=astsa) to get to html help quickly. Another option is to type ?astsa and then use the link to the package Index at the bottom of the page. A pdf version of the astsa manual may be found at CRAN: http: //cran.r-project.org/web/packages/astsa/astsa.pdf. Further details on R packages may be found in Section R.3.2. A useful command to see all the data files available to you, including those loaded with astsa, is data()

## R. 3 Getting Started

The best way to use the rest of this appendix is to start up R and enter the example code as it is presented. Also, you can use the results and help files to get a better understanding of how R works (or doesn't work).

The convention throughout the text is that R code is in blue, output is purple and comments are \# green. Get comfortable, then start her up and try some simple tasks.

```
2+2 # addition
    [1] 5
5*5 + 2 # multiplication and addition
    [1] 27
5/5 - 3 # division and subtraction
    [1] -2
log(exp(pi)) # log, exponential, pi
    [1] 3.141593
sin(pi/2) # sinusoids
    [1] 1
exp(1)^(-2) # power
    [1] 0.1353353
sqrt(8) # square root
    [1] 2.828427
1:5 # sequences
    [1] 1 2 3 4 5
seq(1, 10, by=2) # sequences
    [1] 1 3 5 7 9
rep(2,3) # repeat 2 three times
    [1] 2 2 2
```

Next, we'll use assignment to make some objects:

```
x <- 1 + 2 # put 1 + 2 in object x
x = 1 + 2 # same as above with fewer keystrokes
1 + 2 -> x # same
x # view object x
    [1] 3
    (y = 9*3) # put 9 times 3 in y and view the result
    [1] 27
    (z = rnorm(5)) # put 5 standard normals into z and print z
    [1] 0.96607946 1.98135811 -0.06064527 0.31028473 0.02046853
```

It is worth pointing out R's recycling rule for doing arithmetic. In the code below, $c$ () [concatenation] is used to create a vector. Note the use of the semicolon for multiple commands on one line.

```
x = c(1, 2, 3, 4); y = c(2, 4, 6, 8); z = c(10, 20); w = c(8, 3, 2)
x*y # 1*2, 2*4, 3*6, 4*8
    [1] 2 8 18 32
x+z # 1+10, 2+20, 3+10, 4+20
    [1] 11 22 13 24
y+w # what happened here?
    [1] 10 7 8 16
Warning message:
    In y + w : longer object length is not a multiple of
            shorter object length
```

            To list your objects, remove objects, get help, find out which directory is current
    (or to change it) or to quit, use the following commands:
ls() \# list all objects
"dummy" "mydata" "x" "y" "z"
ls(pattern = "my") \# list every object that contains "my"
"dummy" "mydata"
rm(dummy) \# remove object "dummy"
rm(list=ls()) \# remove almost everything (use with caution)
help.start() \# html help and documentation (use it)
data() \# list of available data sets
help(exp) \# specific help (?exp is the same)
getwd() \# get working directory
setwd() \# change working directory
q() \# end the session (keep reading)

When you quit, R will prompt you to save an image of your current workspace. Answering "yes" will save all the work you have done so far, and load it up when you next start R.
To create your own data set, you can make a data vector as follows:
mydata $=c(1,2,3,2,1)$
Now you have an object called mydata that contains five elements. R calls these objects vectors even though they have no dimensions (no rows, no columns); they do have order and length:

```
mydata # display the data
    [1] 1 2 3 2 1
mydata[3] # the third element
    [1] 3
mydata[3:5] # elements three through five
    [1] 3 2 1
mydata[-(1:2)] # everything except the first two elements
    [1] 3 2 1
length(mydata) # number of elements
    [1] 5
dim(mydata) # no dimensions
    NULL
mydata = as.matrix(mydata) # make it a matrix
dim(mydata) # now it has dimensions
    [1] 5 1
```

If you have an external data set, you can use scan or read.table (or some variant) to input the data. For example, suppose you have an AsciI (text) data file
called dummy.txt in your working directory, and the file looks like this:

```
1223 2 1
90210
(dummy = scan("dummy.txt")) # scan and view it
Read 10 items
    [1] 1 2 3 2 1 9 0 2 10
(dummy = read.table("dummy.txt")) # read and view it
    V1 V2 V3 V4 V5
    1
    9 0 2 1 0
```

There is a difference between scan and read.table. The former produced a data vector of 10 items while the latter produced a data frame with names V1 to V5 and two observations per variate. In this case, if you want to list (or use) the second variate, v2, you would use
dummy\$V2
[1] 20
and so on. You might want to look at the help files ?scan and ?read.table now. Data frames (?data.frame) are "used as the fundamental data structure by most of R's modeling software." Notice that R gave the columns of dummy generic names, V1, ..., V5. You can provide your own names and then use the names to access the data without the use of \$ as above.
colnames(dummy) = c("Dog", "Cat", "Rat", "Pig", "Man")
attach(dummy)
Cat
[1] 20
Rat*(Pig - Man) \# animal arithmetic
[1] 32
detach(dummy) \# clean up (if desired)
$R$ is case sensitive, thus cat and Cat are different. Also, cat is a reserved name (?cat) in R, so using "cat" instead of "Cat" may cause problems later. You may also include a header in the data file to avoid using line 5 above. For example, if the file dummy. txt looks like this,

| Dog | Cat | Rat | Pig | Man |
| :---: | :---: | :---: | :---: | :---: |
| 1 | 2 | 3 | 2 | 1 |
| 9 | 0 | 2 | 1 | 0 |

then use the following command to read the data.
(dummy = read.table("dummy.txt", header=TRUE))

| Dog |  |  |  |  | Cat |
| :---: | ---: | ---: | ---: | ---: | ---: |
| 1 | 1 | 2 | 3 | 2 | 1 |
| 2 | 9 | 0 | 2 | 1 | 0 |

Another popular format for data files is .csv (comma separated values). In this case, use read.csv instead of read.table. The default for .csv files is header=TRUE; type ?read.table for further information.

Two commands that are used frequently to manipulate data are cbind for column binding and rbind for row binding. The following is an example.

```
x = runif(4) # generate 4 values from a uniform(0,1) into object x
y = runif(4) # generate 4 more and put them into object y
(u = cbind(x,y)) # column bind the two vectors (4 by 2 matrix)
```

```
            x y
    [1,] 0.6547304 0.7503984
    [2,] 0.8222048 0.1335557
    [3,] 0.4555755 0.2151735
    [4,] 0.9843289 0.8483795
(u = rbind(x,y)) # row bind the two vectors (2 by 4 matrix)
    [,1] [,2] [,3] [,4]
x 0.6547304 0.8222048 0.4555755 0.9843289
y 0.7503984 0. 1335557 0. 2151735 0.8483795
```


## R.3.1 Basic Statistics

Summary statistics are fairly easy to obtain. We will simulate 25 normals with $\mu=10$ and $\sigma=4$ and then perform some basic analyses.

```
set.seed(90210) # so you can reproduce these results
x = rnorm(25, 10, 4) # generate the data
mean(x)
    [1] 9.473883
median(x) # compute the median
    [1] 9.448511
var(x) # compute the variance
    [1] 13.9267
sd(x) # compute the standard deviation
    [1] 3.73185
max(x) # find the largest value
    [1] 17.32609
which.max(x) # index of the max (x[25] in this case)
    [1] 25
min(x) # find the smallest value
    [1] 2.678173
summary(x) # a five number summary (with the mean added)
            Min. 1st Qu. Median Mean 3rd Qu. Max.
            2.678 7.824 9.449 9.474 11.180 17.330
par(mfrow=c(1,2)) # multifigure setup (?par for info)
hist(x, col="lightblue", prob=TRUE) # see figure below
lines(density(x, bw=2)) # fit a density over it
boxplot(x, main="Boxplot of x", col="lightblue") # see figure below
```



Fig. R.1. Histogram and boxplot of simulated normal data.


Fig. R.2. Crazy example.

We introduced some R graphics without saying much about it. Aside from the manual, there are many useful websites that expand on this topic; e.g., https://www. stat.auckland.ac.nz/~paul/RGraphics/rgraphics.html.

It can't hurt to learn a little about programming in R because you will see some of it along the way. Consider a simple program that we will call crazy to produce a graph of a sequence of sample means of increasing sample sizes from a Cauchy distribution with location parameter zero. The code is:

```
crazy <- function(num) {
    x <- rep(NA, num)
    for (n in 1:num) x[n] <- mean(rcauchy(n))
    plot(x, type="l", xlab="sample size", ylab="sample mean")
    }
```

The first line creates the function crazy and gives it one argument, num, that is the sample size that will end the sequence. Line 2 makes a vector, x , of num missing values NA, that will be used to store the sample means. Line 3 generates $n$ random Cauchy variates [rcauchy ( $n$ )], finds the mean of those values, and puts the result into $\mathrm{x}[\mathrm{n}]$, the $n$-th value of x . The process is repeated in a "do loop" num times so that $x[1]$ is the sample mean from a sample of size one, $x[2]$ is the sample mean from a sample of size two, and so on, until finally, $\mathrm{x}[\mathrm{num}]$ is the sample mean from a sample of size num. After the do loop is complete, the fourth line generates a graphic (see Figure R.2). The fifth line closes the function. To use crazy with a limit sample size of 100 , for example, type
crazy (100)
and you will get a graphic that looks like Figure R. 2

## R.3.2 Packages

We have already discussed how to obtain and load the R package for this text, astsa. You may want to use other R packages as we have done throughout the text. In this case, you have to first download the package and then install it. For example, install.packages(c("wavethresh", "tseries"))
will download and install the packages wavethresh that we use in Chapter 4 and tseries that we use in Chapter 5; you will be asked to choose the closest mirror to you. To use a package, you have to load it at each start up of R, for example:
require(wavethresh) \# load the wavethresh package
A good way to get help for a package is to use html help
help.start()
and follow the Packages link.
To make sure you have the current version of your R packages, periodically update them by issuing the command
update.packages(ask=FALSE)
You will be asked to choose the nearest repository. Using "ask=FALSE" means you won't be asked about updating each package. If you want a list of your packages, run installed.packages() [, 1:3]
giving the first 3 columns of info: the [1] package name, [2] library path, and [3] version number.

## R.3.3 Word

Finally, a word of caution: TRUE and FALSE are reserved words, whereas T and F are initially set to these. Get in the habit of using the words rather than the letters T or F because you may get into trouble if you do something like $\mathrm{F}=\mathrm{qf}(\mathrm{p}=.01, \mathrm{df} 1=3, \mathrm{df} 2=9)$
so that F is no longer FALSE, but a quantile of the specified $F$-distribution. Better yet, just remember not to use T or F for anything.

## R. 4 Time Series Primer

In this section, we give a brief introduction on using R for time series. We assume that astsa has been loaded. To create a time series object, use the command ts. Related commands are as.ts to coerce an object to a time series and is.ts to test whether an object is a time series.

First, make a small data set:
(mydata $=\mathrm{c}(1,2,3,2,1)$ ) \# make it and view it
[1] 12321
Now make it a time series:
(mydata $=$ as.ts(mydata))
Time Series:
Start = 1
End $=5$
Frequency $=1$
[1] 12321
Make it an annual time series that starts in 1950:
(mydata $=$ ts (mydata, start=1950))
Time Series:
Start = 1950
End = 1954

```
Frequency = 1
```

[1] 12321
Now make it a quarterly time series that starts in 1950-III:

```
(mydata = ts(mydata, start=c(1950,3), frequency=4))
            Qtr1 Qtr2 Qtr3 Qtr4
    1950 1 2
    1951 3 2 1
time(mydata) # view the sampled times
    Qtr1 Qtr2 Qtr3 Qtr4
    1950 1950.50 1950.75
    1951 1951.00 1951.25 1951.50
```

To use part of a time series object, use window ():

```
(x = window(mydata, start=c(1951,1), end=c(1951,3)))
        Qtr1 Qtr2 Qtr3
    1951 3 2 1
```

Next, we'll look at lagging and differencing. First make a simple series, $x_{t}$ :
$\mathrm{x}=\mathrm{ts}(1: 5)$
Now, column bind (cbind) lagged values of $x_{t}$ and you will notice that $\operatorname{lag}(\mathrm{x})$ is forward lag, whereas lag $(\mathrm{x},-1)$ is backward lag.
cbind(x, $\operatorname{lag}(x), \operatorname{lag}(x,-1))$
$x \operatorname{lag}(x) \operatorname{lag}(x,-1)$

| 0 | $N A$ | 1 | $N A$ |  |
| :--- | ---: | ---: | ---: | :--- |
| 1 | 1 | 2 | $N A$ |  |
| 2 | 2 | 3 | 1 | <- in this row, for example, $x$ is 3, |
| 3 | 3 | 4 | 3 | $\operatorname{lag}(x)$ is ahead at 4, and |
| 4 | 4 | 5 | 4 | $\operatorname{lag}(x,-1)$ is behind at 2 |
| 5 | 5 | $N A$ | 5 |  |
| 6 | $N A$ | $N A$ |  |  |

Compare cbind and ts.intersect:
ts.intersect(x, $\operatorname{lag}(x, 1), \operatorname{lag}(x,-1))$
Time Series: Start $=2$ End $=4$ Frequency $=1$
$x \operatorname{lag}(x, 1) \operatorname{lag}(x,-1)$
$\begin{array}{llll}2 & 2 & 3 & 1\end{array}$
$\begin{array}{llll}3 & 3 & 4 & 2\end{array}$
$\begin{array}{llll}4 & 4 & 5 & 3\end{array}$
To difference a series, $\nabla x_{t}=x_{t}-x_{t-1}$, use
$\operatorname{diff}(x)$
but note that
$\operatorname{diff}(x, 2)$
is not second order differencing, it is $x_{t}-x_{t-2}$. For second order differencing, that is, $\nabla^{2} x_{t}$, do one of these:
diff(diff(x))
diff(x, diff=2) \# same thing
and so on for higher order differencing.
For graphing time series, there a few standard plotting mechanisms that we use repeatedly. If $x$ is a time series, then $p l o t(x)$ will produce a time plot. If $x$ is not a time series object, then plot.ts( x ) will coerce it into a time plot as will ts.plot(x). There are differences, which we explore in the following. It would be a good idea to skim the graphical parameters help file (?par) while you are here. See Figure R. 3 for the resulting graphic.


Fig. R.3. Demonstration of different R graphic tools for plotting time series.

```
x = -5:5 # x is NOT a time series object
y = 5*}\operatorname{cos}(x) # neither is y
par(mfrow=c(3,2)) # multifigure setup: 3 rows, 2 cols
plot(x, main="plot(x)")
plot(x, y, main="plot(x,y)")
plot.ts(x, main="plot.ts(x)")
plot.ts(x, y, main="plot.ts(x,y)")
ts.plot(x, main="ts.plot(x)")
ts.plot(ts(x), ts(y), col=1:2, main="ts.plot(x,y)")
```

We will also make use of regression via $\operatorname{lm}()$. First, suppose we want to fit a simple linear regression, $y=\alpha+\beta x+\epsilon$. In R , the formula is written as $\mathrm{y} \sim \mathrm{x}$ :

```
set.seed(1999) # so you can reproduce the result
```

$\mathrm{x}=\operatorname{rnorm}(10,0,1)$
$y=x+\operatorname{rnorm}(10,0,1)$
summary (fit <- $\operatorname{lm}(y \sim x)$ )

| Residuals: |  |  |  |  |
| :--- | ---: | ---: | ---: | ---: |
| Min | $1 Q$ | Median | $3 Q$ | Max |
| -0.8851 | -0.3867 | 0.1325 | 0.3896 | 0.6561 |

Coefficients:
Estimate Std. Error t value $\operatorname{Pr}(>\mathrm{t})$

| (Intercept) | 0.2576 | 0.1892 | 1.362 | 0.2104 |
| :--- | :--- | :--- | :--- | :--- |
| x | 0.4577 | 0.2016 | 2.270 | 0.0529 |

```
    Residual standard error: 0.58 on 8 degrees of freedom
    Multiple R-squared: 0.3918, Adjusted R-squared: 0.3157
    F-statistic: 5.153 on 1 and 8 DF, p-value: 0.05289
plot(x, y) # draw a scatterplot of the data (not shown)
abline(fit) # add the fitted line to the plot (not shown)
```

All sorts of information can be extracted from the lm object, which we called fit. For example,

```
resid(fit) # will display the residuals (not shown)
fitted(fit) # will display the fitted values (not shown)
lm(y ~ 0 + x) # will exclude the intercept (not shown)
```

You have to be careful if you use $\operatorname{lm}()$ for lagged values of a time series. If you use $\operatorname{lm}()$, then what you have to do is "tie" the series together using ts.intersect. If you do not tie the series together, they will not be aligned properly. Please read the warning Using time series in the $\operatorname{lm}()$ help file [help(lm)]. Here is an example regressing astsa data, weekly cardiovascular mortality (cmort) on particulate pollution (part) at the present value and lagged four weeks (part4). First, we create ded, which consists of the intersection of the three series:
ded $=$ ts.intersect(cmort, part, part4=lag(part,-4))
Now the series are all aligned and the regression will work.

| Coefficients: |  |  |  |
| :---: | :---: | :---: | :---: |
|  | Estimate | Std. Error | t value |
| (Intercept) | 69.01020 | 1.37498 | 50.190 |
| part | 0.15140 | 0.02898 | 5.225 |
| part4 | 0.26297 | 0.02899 | 9.071 |

Residual standard error: 8.323 on 501 degrees of freedom Multiple R-Squared: 0.3091, Adjusted R-squared: 0.3063 F-statistic: 112.1 on 2 and $501 \mathrm{DF}, \mathrm{p}$-value: < $2.2 \mathrm{e}-16$
There was no need to rename lag (part, -4 ) to part4, it's just an example of what you can do. There is a package called dynlm that makes it easy to fit lagged regressions. The basic advantage of dynlm is that it avoids having to make a data frame; that is, line 2 would be avoided.

In Problem 2.1, you are asked to fit a regression model

$$
x_{t}=\beta t+\alpha_{1} Q_{1}(t)+\alpha_{2} Q_{2}(t)+\alpha_{3} Q_{3}(t)+\alpha_{4} Q_{4}(t)+w_{t}
$$

where $x_{t}$ is logged Johnson \& Johnson quarterly earnings ( $n=84$ ), and $Q_{i}(t)$ is the indicator of quarter $i=1,2,3,4$. The indicators can be made using factor.

```
trend = time(jj) - 1970 # helps to 'center' time
Q = factor(cycle(jj)) # make (Q)uarter factors
reg = lm(log(jj)~0 + trend + Q, na.action=NULL) # no intercept
model.matrix(reg) # view the model matrix
        trend Q1 Q2 Q3 Q4
    1 -10.00 1 0 0 0
    2 -9.75 00 1 0 0
    -9.50 0 0 1 0
```

```
    -9.25 0 0 0 1
    . . . . .
    10.50 0 0 1 0
    10.75 0 0 0 1
summary(reg) # view the results (not shown)
```

The workhorse for ARIMA simulations is arima.sim. Here are some examples; no output is shown here so you're on your own.

```
\(\mathrm{x}=\operatorname{arima} . \operatorname{sim}(\operatorname{list}(\operatorname{order}=\mathrm{c}(1,0,0), \operatorname{ar}=.9), \mathrm{n}=100)+50\) \# AR(1) w/mean 50
\(\mathrm{x}=\operatorname{arima} . \operatorname{sim}(\operatorname{list}(\operatorname{order}=\mathrm{c}(2,0,0), \operatorname{ar}=\mathrm{c}(1,-.9)), \mathrm{n}=100) \quad \# \operatorname{AR}(2)\)
\(\mathrm{x}=\operatorname{arima} . \operatorname{sim}(\operatorname{list}(\operatorname{order}=\mathrm{c}(1,1,1)\), ar=.9, ma=-.5) , \(\mathrm{n}=200)\) \# ARIMA(1,1,1)
```

Next, we'll discuss ARIMA estimation. This gets a bit tricky because R is not useR friendly when it comes to fitting ARIMA models. Much of the story is spelled out in the "R Issues" page of the website for the text. In Chapter 3, we use the scripts acf2, sarima, and sarima.for that are included with astsa. But we will also show you how to use the scripts included with $R$.

First, we'll fit an $\operatorname{ARMA}(1,1)$ model to some simulated data (with diagnostics and forecasting):

```
set.seed(666)
x = 50 + arima.sim(list(order=c(1,0,1), ar=.9, ma=-.5), n=200)
acf(x); pacf(x) # display sample ACF and PACF ... or ...
acf2(x) # use our script (no output shown)
(x.fit = arima(x, order = c(1, 0, 1))) # fit the model
    Call: arima(x = x, order = c(1, 0, 1))
    Coefficients:
        ar1 ma1 intercept
        0.8340-0.432 49.8960
    s.e. 0.0645 0.111 0.2452
    sigma^2 estimated as 1.070: log likelihood = -290.79, aic = 589.58
```

Note that the reported intercept estimate is an estimate of the mean and not the constant. That is, the fitted model is

$$
\widehat{x}_{t}-49.896=.834\left(x_{t-1}-49.896\right)+\widehat{w}_{t}
$$

where $\widehat{\sigma}_{w}^{2}=1.070$. Incorrect diagnostics can be accomplished as follows:
tsdiag(x.fit, gof.lag=20) \# !!!! don't use this !!!!
That's right, the Ljung-Box-Pierce test is not correct because it does not take into account the fact that the residuals are from a fitted model. If the analysis is repeated using the sarima script, a partial output would look like the following (sarima will also display the correct diagnostics as a graphic; e.g., see Figure 3.17 and Figure R.4): sarima(x, 1, 0, 1)

Coefficients:

```
                    ar1 ma1 xmean
            0.8340-0.432 49.8960
    s.e. 0.0645 0.111 0.2452
    sigma^2 estimated as 1.070: log likelihood = -290.79, aic = 589.58
    $AIC [1] 1.097494 $AICc [1] 1.108519 $BIC [1] 0.1469684
```

An easy way to see a summary of the significance tests is to use the R package lmtest, which has to be downloaded first:
require('lmtest') \# addicted to p-values?
coeftest(sarima(x, 1,0,1)\$fit)


Fig. R.4. Diagnostics from sarima

|  | Estimate | Std.Error | z.value $\quad$ p.value |  |
| :--- | ---: | ---: | ---: | ---: |
| ar1 | 0.83405 | 0.06448 | $12.9350<2.2 \mathrm{e}-16$ |  |
| ma1 | -0.43199 | 0.11104 | -3.8904 | 0.0001001 |
| xmean | 49.89622 | 0.24519 | $203.5003<2.2 e-16$ |  |

Note that the results of the sarima fit are in fit. Forecasting using the script sarima.for can be accomplished in one line:
sarima.for(x, 10, 1, 0, 1)
If you don't use astsa, you would obtain and plot the forecasts using the following:
$\mathrm{x} . \mathrm{fore}^{=}$predict(x.fit, n.ahead=10)
$\mathrm{U}=\mathrm{x}$. fore\$pred $+2 * \mathrm{x}$.fore\$se \# x.fore\$pred holds predicted values
$\mathrm{L}=\mathrm{x} . f 0 r e \$ p r e d-2 * x$.fore\$se \# x.fore\$se holds stnd errors
$\min y=\min (x, L) ; \quad \operatorname{maxy}=\max (x, U)$
ts.plot(x, x.fore\$pred, col=1:2, ylim=c(miny, maxy))
lines(U, col="blue", lty="dashed")
lines(L, col="blue", lty="dashed")
We close this appendix with a quick spectral analysis. This material is covered in detail in Chapter 4, so we will not discuss this example in much detail here. We will simulate an $\operatorname{AR}(2)$ and then estimate the spectrum via nonparametric and parametric methods. No graphics are shown, but we have confidence that you are proficient enough in R to display them yourself.

```
x = arima.sim(list(order=c(2,0,0), ar=c(1,-.9)), n=2^8) # some data
(u = polyroot(c(1,-1,.9))) # x is AR(2) w/complex roots
    [1] 0.5555556+Q.8958064i 0.5555556-0.8958064i
Arg(u[1])/(2*pi) # dominant frequency around . 16
```

```
    [1] 0.1616497
par(mfcol=c(2,2))
plot.ts(x)
mvspec(x, spans=c(5,5), taper=.1, log="no") # nonparametric estimate
spec.ar(x, log="no") # parametric spectral estimate
arma.spec(ar=c(1,-.9), log="no") # true spectral density
The script arma.spec is included in astsa. See spectrum and spec.pgram as alter-
natives to mvspec, which is part of astsa. Finally, note the easiest way to get a raw
periodogram is:
per = abs(fft(x))^2/length(x) # abs() and Mod() are same here
```


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[^0]:    ** Throughout the text, we assume that the R package for the book, astsa, has been downloaded and installed. See Appendix R (Section R.2) for further details.

[^1]:    ${ }^{1}$ Expectation is discussed in the third chapter of the basic probability facts pdf mentioned in the preface. For continuous-valued finite variance processes, the mean is $\mu_{t}=E\left(x_{t}\right)=\int_{-\infty}^{\infty} x f_{t}(x) d x$ and the variance is $\sigma_{t}^{2}=E\left(x_{t}-\mu_{t}\right)^{2}=\int_{-\infty}^{\infty}\left(x-\mu_{t}\right)^{2} f_{t}(x) d x$, where $f_{t}$ is the density of $x_{t}$. If $x_{t}$ is Gaussian with mean $\mu_{t}$ and variance $\sigma_{t}^{2}$, abbreviated as $x_{t} \sim \mathrm{~N}\left(\mu_{t}, \sigma_{t}^{2}\right)$, the marginal density is given by $f_{t}(x)=\frac{1}{\sigma_{t} \sqrt{2 \pi}} \exp \left\{-\frac{1}{2 \sigma_{t}^{2}}\left(x-\mu_{t}\right)^{2}\right\}$ for $x \in \mathbb{R}$.

[^2]:    ${ }^{2}$ The Cauchy-Schwarz inequality implies $|\gamma(s, t)|^{2} \leq \gamma(s, s) \gamma(t, t)$.

[^3]:    ${ }^{3}$ In this text, $z_{.025}=1.95996398454005423552 \ldots$ of normal fame, which is often rounded to 1.96 , is rounded to 2 .

[^4]:    ${ }^{1}$ Formally, AIC is defined as $-2 \log L_{k}+2 k$ where $L_{k}$ is the maximum value of the likelihood and $k$ is the number of parameters in the model. For the normal regression problem, AIC can be reduced to the form given by (2.15).

[^5]:    ${ }^{2}$ Because the error term, $y_{t}$, is not assumed to be iid, the reader may feel that weighted least squares is called for in this case. The problem is, we do not know the behavior of $y_{t}$ and that is precisely what we are trying to assess at this stage. An important result by Grenander and Rosenblatt (1957, Ch 7), however, is that under mild conditions, for polynomial or periodic regression, when the sample size is large, ordinary least squares is equivalent to weighted least squares in that the standard errors of the estimates are the approximately same.

[^6]:    ${ }^{3} \cos (\alpha \pm \beta)=\cos (\alpha) \cos (\beta) \mp \sin (\alpha) \sin (\beta)$.

[^7]:    ${ }^{1}$ Iterate backward, $x_{t}=\phi x_{t-1}+w_{t}=\phi\left(\phi x_{t-2}+w_{t-1}\right)+w_{t}=\phi^{2} x_{t-2}+\phi w_{t-1}+w_{t}=\cdots=$ $\phi^{k} x_{t-k}+\sum_{j=0}^{k-1} \phi^{j} w_{t-j}$. If $|\phi|<1$ and $\sup _{t} E\left(x_{t}^{2}\right)<\infty$, then $\lim _{k \rightarrow \infty} E\left(x_{t}-\sum_{j=0}^{k-1} \phi^{j} w_{t-j}\right)^{2}=$ $\lim _{k \rightarrow \infty} \phi^{2 k} E\left(x_{t-k}^{2}\right)=0$.

[^8]:    ${ }^{2}$ Some texts and software packages write the MA model with negative coefficients; that is, $x_{t}=$ $w_{t}-\theta_{1} w_{t-1}-\theta_{2} w_{t-2}-\cdots-\theta_{q} w_{t-q}$.

[^9]:    * $\phi(z)$ can't be zero in here. . you wouldn't want to divide by zero, would you?
    $\dagger \theta(z)$ can't be zero in here.

[^10]:    ${ }^{3}$ Newton's method and Taylor expansion (links to WikiBooks K-12 calculus book).

[^11]:    ${ }^{4}$ To estimate the standard error, we are using the standard regression results from (2.6) as an approximation

[^12]:    ${ }^{5} X_{n} \sim \mathrm{AN}\left(\mu_{n}, \sigma_{n}^{2}\right)$ if $\operatorname{Pr}\left\{\left(X_{n}-\mu_{n}\right) / \sigma_{n} \leq z\right\} \rightarrow \operatorname{Pr}\{Z \leq z\}$ as $n \rightarrow \infty$, where $Z \sim \mathrm{~N}(0,1)$.

[^13]:    ${ }^{6}$ To establish (3.51), by invertibility, the $w_{n}, w_{n-1}, \ldots$ are functions of the observations $x_{n}, x_{n-1}, \ldots$. Letting $w_{n+k}^{n}=E\left(w_{n+k} \mid x_{n}, x_{n-1}, \ldots\right)$, we have $w_{n+k}^{n}=w_{n+k}$ for $k \leq 0$, but for future errors, $w_{n+k}^{n}=E w_{n+k}=0$ for $k \geq 1$. Thus

    $$
    x_{n+m}=\sum_{j=0}^{\infty} \psi_{j} w_{m+n-j} \quad \text { and } \quad x_{n+m}^{n}=\sum_{j=0}^{\infty} \psi_{j} w_{m+n-j}^{n}=\sum_{j=m}^{\infty} \psi_{j} w_{m+n-j}
    $$

    so that $E\left[x_{n+m}-x_{n+m}^{n}\right]^{2}=E\left[\sum_{j=0}^{m-1} \psi_{j} w_{n+m-j}\right]^{2}$, which is (3.51).

[^14]:    ** This section may be skipped without harming any living things.

[^15]:    ${ }^{7}$ The criterion function is sometimes called the profile or concentrated likelihood.

[^16]:    ${ }^{8} \log (1+p)=p-\frac{p^{2}}{2}+\frac{p^{3}}{3}-\cdots$ for $-1<p \leq 1$. If $p$ is a small percent-change, then the higher-order terms in the expansion are negligible.

[^17]:    ${ }^{9}$ If you're addicted to p-values, see page 187 in Appendix R.

[^18]:    ${ }^{10}$ Write the model in vector notation as $\boldsymbol{y}=Z \boldsymbol{\beta}+\boldsymbol{x}$, and let $\Gamma=\left\{\gamma_{x}(s, t)\right\}$. Then $\Gamma^{-1 / 2} \boldsymbol{y}=\Gamma^{-1 / 2} Z \boldsymbol{\beta}+$ $\Gamma^{-1 / 2} \boldsymbol{x}$, so that we can write the model as $\boldsymbol{y}^{*}=Z^{*} \boldsymbol{\beta}+\boldsymbol{\delta}$, where the covariance matrix of $\boldsymbol{\delta}$ is the identity. It follows that the weighted estimate of $\boldsymbol{\beta}$ is $\widehat{\boldsymbol{\beta}}_{w}=\left(Z^{\prime} \Gamma^{-1} Z\right)^{-1} Z^{\prime} \Gamma^{-1} \boldsymbol{y}$, and the variancecovariance matrix of the estimator is $\operatorname{var}\left(\widehat{\boldsymbol{\beta}}_{w}\right)=\left(Z^{\prime} \Gamma^{-1} Z\right)^{-1}$. If $x_{t}$ is white noise, then $\Gamma=\sigma^{2} I$ and these results reduce to the usual least squares results.

[^19]:    ${ }^{1} \cos (\alpha \pm \beta)=\cos (\alpha) \cos (\beta) \mp \sin (\alpha) \sin (\beta)$.

[^20]:    ${ }^{2}$ Useful information: Euler's formula: $e^{i \alpha}=\cos (\alpha)+i \sin (\alpha)$. Consequently, $\cos (\alpha)=\frac{e^{i \alpha}+e^{-i \alpha}}{2}$, and $\sin (\alpha)=\frac{e^{i \alpha}-e^{-i \alpha}}{2 i}$. Also, $\frac{1}{i}=-i$ because $-i \times i=1$. If $z=a+i b$ is complex, then $|z|^{2}=z \bar{z}=$ $(a+i b)(a-i b)=a^{2}+b^{2}$.

[^21]:    ${ }^{3}$ Recall that if $z=a+i b$, then $\bar{z}=a-i b$, and $|z|^{2}=z \bar{z}=a^{2}+b^{2}$.

[^22]:    ${ }^{4}$ The DFTs of $x_{t}$ and of ( $x_{t}-\bar{x}$ ) are the same except at the zero frequency. This follows because $\sum_{t=1}^{n} \exp \left(-2 \pi i t \frac{j}{n}\right)=0$. Consequently, $\left|d\left(\omega_{j}\right)\right|^{2}=n^{-1} \sum_{t=1}^{n} \sum_{s=1}^{n}\left(x_{t}-\bar{x}\right)\left(x_{s}-\bar{x}\right) e^{-2 \pi i \omega_{j}(t-s)}=$ $n^{-1} \sum_{h=-(n-1)}^{n-1} \sum_{t=1}^{n-|h|}\left(x_{t+|h|}-\bar{x}\right)\left(x_{t}-\bar{x}\right) e^{-2 \pi i \omega_{j} h}$, which is (4.20).

[^23]:    ${ }^{5}$ For this section, it will be useful to recall the facts $e^{-i \alpha}=\cos (\alpha)-i \sin (\alpha)$ and if $z=a+i b$, then $\bar{z}=a-i b$.

[^24]:    ${ }^{6}$ If $Z$ is a complex matrix, then $Z^{*}=\bar{Z}^{\prime}$ denotes the conjugate transpose operation. That is, $Z^{*}$ is the result of replacing each element of $Z$ by its complex conjugate and transposing the resulting matrix.

[^25]:    ${ }^{7}$ If $L=1$ then $\bar{\rho}_{y \cdot x}^{2}(\omega) \equiv 1$.

