## Contents

## Lecture Notes for Econometrics 2002 (first year

PhD course in Stockholm)

## Paul Söderlind

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${ }^{1}$ University of St. Gallen and CEPR. Address: s/bf-HSG, Rosenbergstrasse 52, CH-9000 St. Gallen, Switzerland. E-mail: Paul.Soderlind@unisg.ch. Document name: EcmAll.TeX.
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## 1 Introduction

### 1.1 Least Squares

Consider the simplest linear model

$$
\begin{equation*}
y_{t}=x_{t} \beta_{0}+u_{t} \tag{1.1}
\end{equation*}
$$

where all variables are zero mean scalars and where $\beta_{0}$ is the true value of the parameter we want to estimate. The task is to use a sample $\left\{y_{t}, x_{t}\right\}_{t=1}^{T}$ to estimate $\beta$ and to test hypotheses about its value, for instance that $\beta=0$.

If there were no movements in the unobserved errors, $u_{t}$, in (1.1), then any sample would provide us with a perfect estimate of $\beta$. With errors, any estimate of $\beta$ will still leave us with some uncertainty about what the true value is. The two perhaps most important issues is econometrics are how to construct a good estimator of $\beta$ and how to assess the uncertainty about the true value.

For any possible estimate, $\hat{\beta}$, we get a fitted residual

$$
\begin{equation*}
\hat{u}_{t}=y_{t}-x_{t} \hat{\beta} \tag{1.2}
\end{equation*}
$$

One appealing method of choosing $\hat{\beta}$ is to minimize the part of the movements in $y_{t}$ that we cannot explain by $x_{t} \hat{\beta}$, that is, to minimize the movements in $\hat{u}_{t}$. There are several candidates for how to measure the "movements," but the most common is by the mean of squared errors, that is, $\Sigma_{t=1}^{T} \hat{u}_{t}^{2} / T$. We will later look at estimators where we instead use $\Sigma_{t=1}^{T}\left|\hat{u}_{t}\right| / T$.

With the sum or mean of squared errors as the loss function, the optimization problem

$$
\begin{equation*}
\min _{\beta} \frac{1}{T} \sum_{t=1}^{T}\left(y_{t}-x_{t} \beta\right)^{2} \tag{1.3}
\end{equation*}
$$

has the first order condition that the derivative should be zero as the optimal estimate $\hat{\beta}$

$$
\begin{equation*}
\frac{1}{T} \sum_{t=1}^{T} x_{t}\left(y_{t}-x_{t} \hat{\beta}\right)=0 \tag{1.4}
\end{equation*}
$$

which we can solve for $\hat{\beta}$ as

$$
\begin{align*}
\hat{\beta} & =\left(\frac{1}{T} \sum_{t=1}^{T} x_{t}^{2}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} x_{t} y_{t}, \text { or }  \tag{1.5}\\
& =\widehat{\operatorname{Var}}\left(x_{t}\right)^{-1} \widehat{\operatorname{Cov}}\left(x_{t}, y_{t}\right) \tag{1.6}
\end{align*}
$$

where a hat indicates a sample estimate. This is the Least Squares (LS) estimator.

### 1.2 Maximum Likelihood

A different route to arrive at an estimator is to maximize the likelihood function. If $u_{t}$ in (1.1) is iid $N\left(0, \sigma^{2}\right)$, then the probability density function of $u_{t}$ is

$$
\begin{equation*}
\operatorname{pdf}\left(u_{t}\right)=\sqrt{2 \pi \sigma^{2}} \exp \left[-u_{t}^{2} /\left(2 \sigma^{2}\right)\right] \tag{1.7}
\end{equation*}
$$

Since the errors are independent, we get the joint pdf of the $u_{1}, u_{2}, \ldots, u_{T}$ by multiplying the marginal pdfs of each of the errors. Then substitute $y_{t}-x_{t} \beta$ for $u_{t}$ (the derivative of the transformation is unity) and take logs to get the $\log$ likelihood function of the sample

$$
\begin{equation*}
\ln L=-\frac{T}{2} \ln (2 \pi)-\frac{T}{2} \ln \left(\sigma^{2}\right)-\frac{1}{2} \sum_{t=1}^{T}\left(y_{t}-x_{t} \beta\right)^{2} / \sigma^{2} \tag{1.8}
\end{equation*}
$$

This likelihood function is maximized by minimizing the last term, which is proportional to the sum of squared errors - just like in (1.3): LS is ML when the errors are iid normally distributed.

Maximum likelihood estimators have very nice properties, provided the basic distributional assumptions are correct. If they are, then MLE are typically the most efficient/precise estimators, at least asymptotically. ML also provides a coherent framework for testing hypotheses (including the Wald, LM, and LR tests).

### 1.3 The Distribution of $\hat{\beta}$

Equation (1.5) will give different values of $\hat{\beta}$ when we use different samples, that is different draws of the random variables $u_{t}, x_{t}$, and $y_{t}$. Since the true value, $\beta_{0}$, is a fixed constant, this distribution describes the uncertainty we should have about the true value after having obtained a specific estimated value.

To understand the distribution of $\hat{\beta}$, use (1.1) in (1.5) to substitute for $y_{t}$

$$
\begin{align*}
\hat{\beta} & =\left(\frac{1}{T} \sum_{t=1}^{T} x_{t}^{2}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} x_{t}\left(x_{t} \beta_{0}+u_{t}\right) \\
& =\beta_{0}+\left(\frac{1}{T} \sum_{t=1}^{T} x_{t}^{2}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} x_{t} u_{t} \tag{1.9}
\end{align*}
$$

where $\beta_{0}$ is the true value.
The first conclusion from (1.9) is that, with $u_{t}=0$ the estimate would always be perfect - and with large movements in $u_{t}$ we will see large movements in $\hat{\beta}$. The second conclusion is that not even a strong opinion about the distribution of $u_{t}$, for instance that $u_{t}$ is iid $N\left(0, \sigma^{2}\right)$, is enough to tell us the whole story about the distribution of $\hat{\beta}$. The reason is that deviations of $\hat{\beta}$ from $\beta_{0}$ are a function of $x_{t} u_{t}$, not just of $u_{t}$. Of course, when $x_{t}$ are a set of deterministic variables which will always be the same irrespective of which sample we use, then $\hat{\beta}-\beta_{0}$ is a time invariant linear function of $u_{t}$, so the distribution of $u_{t}$ carries over to the distribution of $\hat{\beta}$. This is probably an unrealistic case, which forces us to look elsewhere to understand the properties of $\hat{\beta}$.

There are two main routes to learn more about the distribution of $\hat{\beta}$ : (i) set up a small "experiment" in the computer and simulate the distribution or (ii) use the asymptotic distribution as an approximation. The asymptotic distribution can often be derived, in contrast to the exact distribution in a sample of a given size. If the actual sample is large, then the asymptotic distribution may be a good approximation.

A law of large numbers would (in most cases) say that both $\sum_{t=1}^{T} x_{t}^{2} / T$ and $\sum_{t=1}^{T} x_{t} u_{t} / T$ in (1.9) converges to their expected values as $T \rightarrow \infty$. The reason is that both are sample averages of random variables (clearly, both $x_{t}^{2}$ and $x_{t} u_{t}$ are random variables). These expected values are $\operatorname{Var}\left(x_{t}\right)$ and $\operatorname{Cov}\left(x_{t}, u_{t}\right)$, respectively (recall both $x_{t}$ and $u_{t}$ have zero means). The key to show that $\hat{\beta}$ is consistent, that is, has a probability limit equal to $\beta_{0}$,
is that $\operatorname{Cov}\left(x_{t}, u_{t}\right)=0$. This highlights the importance of using good theory to derive not only the systematic part of (1.1), but also in understanding the properties of the errors. For instance, when theory tells us that $y_{t}$ and $x_{t}$ affect each other (as prices and quantities typically do), then the errors are likely to be correlated with the regressors - and LS is inconsistent. One common way to get around that is to use an instrumental variables technique. More about that later. Consistency is a feature we want from most estimators, since it says that we would at least get it right if we had enough data.

Suppose that $\hat{\beta}$ is consistent. Can we say anything more about the asymptotic distribution. Well, the distribution of $\hat{\beta}$ converges to a spike with all the mass at $\beta_{0}$, but the distribution of $\sqrt{T} \hat{\beta}$, or $\sqrt{T}\left(\hat{\beta}-\beta_{0}\right)$, will typically converge to a non-trivial normal distribution. To see why, note from (1.9) that we can write

$$
\begin{equation*}
\sqrt{T}\left(\hat{\beta}-\beta_{0}\right)=\left(\frac{1}{T} \sum_{t=1}^{T} x_{t}^{2}\right)^{-1} \frac{\sqrt{T}}{T} \sum_{t=1}^{T} x_{t} u_{t} \tag{1.10}
\end{equation*}
$$

The first term on the right hand side will typically converge to the inverse of $\operatorname{Var}\left(x_{t}\right)$, as discussed earlier. The second term is $\sqrt{T}$ times a sample average (of the random variable $x_{t} u_{t}$ ) with a zero expected value, since we assumed that $\hat{\beta}$ is consistent. Under weak conditions, a central limit theorem applies so $\sqrt{T}$ times a sample average converges to a normal distribution. This shows that $\sqrt{T} \hat{\beta}$ has an asymptotic normal distribution. It turns out that this is a property of many estimators, basically because most estimators are some kind of sample average. For an example of a central limit theorem in action, see Appendix B

### 1.4 Diagnostic Tests

Exactly what the variance of $\sqrt{T}\left(\hat{\beta}-\beta_{0}\right)$ is, and how it should be estimated, depends mostly on the properties of the errors. This is one of the main reasons for diagnostic tests. The most common tests are for homoskedastic errors (equal variances of $u_{t}$ and $u_{t-s}$ ) and no autocorrelation (no correlation of $u_{t}$ and $u_{t-s}$ ).

When ML is used, it is common to investigate if the fitted errors satisfy the basic assumptions, for instance, of normality.


### 1.5 Testing Hypotheses about $\hat{\beta}$

Suppose we now that the asymptotic distribution of $\hat{\beta}$ is such that

$$
\begin{equation*}
\sqrt{T}\left(\hat{\beta}-\beta_{0}\right) \xrightarrow{d} N\left(0, v^{2}\right) \text { or } \tag{1.11}
\end{equation*}
$$

We could then test hypotheses about $\hat{\beta}$ as for any other random variable. For instance, consider the hypothesis that $\beta_{0}=0$. If this is true, then

$$
\begin{equation*}
\operatorname{Pr}(\sqrt{T} \hat{\beta} / v<-2)=\operatorname{Pr}(\sqrt{T} \hat{\beta} / v>2) \approx 0.025 \tag{1.12}
\end{equation*}
$$

which says that there is only a $2.5 \%$ chance that a random sample will deliver a value of $\sqrt{T} \hat{\beta} / v$ less than -2 and also a $2.5 \%$ chance that a sample delivers a value larger than 2 , assuming the true value is zero.

We then say that we reject the hypothesis that $\beta_{0}=0$ at the $5 \%$ significance level ( $95 \%$ confidence level) if the test statistics $|\sqrt{T} \hat{\beta} / v|$ is larger than 2 . The idea is that, if the hypothesis is true $\left(\beta_{0}=0\right)$, then this decision rule gives the wrong decision in $5 \%$ of the cases. That is, $5 \%$ of all possible random samples will make us reject a true hypothesis. Note, however, that since this test can only be taken to be an approximation since it relies on the asymptotic distribution, which is an approximation of the true (and typically unknown) distribution.

The natural interpretation of a really large test statistics, $|\sqrt{T} \hat{\beta} / v|=3$ say, is that it is very unlikely that this sample could have been drawn from a distribution where the hypothesis $\beta_{0}=0$ is true. We therefore choose to reject the hypothesis. We also hope that the decision rule we use will indeed make us reject false hypothesis more often than we
reject true hypothesis. For instance, we want the decision rule discussed above to reject $\beta_{0}=0$ more often when $\beta_{0}=1$ than when $\beta_{0}=0$.

There is clearly nothing sacred about the $5 \%$ significance level. It is just a matter of convention that the $5 \%$ and $10 \%$ are the most widely used. However, it is not uncommon to use the $1 \%$ or the $20 \%$. Clearly, the lower the significance level, the harder it is to reject a null hypothesis. At the $1 \%$ level it often turns out that almost no reasonable hypothesis can be rejected.

The t-test described above works only the null hypothesis contains a single restriction. We have to use another approach whenever we want to test several restrictions jointly. The perhaps most common approach is a Wald test. To illustrate the idea, suppose $\beta$ is an $m \times 1$ vector and that $\sqrt{T} \hat{\beta} \xrightarrow{d} N(\mathbf{0}, V)$ under the null hypothesis, where $V$ is a covariance matrix. We then know that

$$
\begin{equation*}
\sqrt{T} \hat{\beta}^{\prime} V^{-1} \sqrt{T} \hat{\beta} \xrightarrow{d} \chi^{2}(m) \tag{1.13}
\end{equation*}
$$

The decision rule is then that if the left hand side of (1.13) is larger that the $5 \%$, say, critical value of the $\chi^{2}(m)$ distribution, then we reject the hypothesis that all elements in $\beta$ are zero.

## A Practical Matters

## A.0. 1 Software

- Gauss, MatLab, RATS, Eviews, Stata, PC-Give, Micro-Fit, TSP, SAS
- Software reviews in The Economic Journal and Journal of Applied Econometrics


## A. 0.2 Useful Econometrics Literature

1. Greene (2000), Econometric Analysis (general)
2. Hayashi (2000), Econometrics (general)
3. Johnston and DiNardo (1997), Econometric Methods (general, fairly easy)
4. Pindyck and Rubinfeld (1997), Econometric Models and Economic Forecasts (general, easy)
5. Verbeek (2000), A Guide to Modern Econometrics (general, easy, good applications)
6. Davidson and MacKinnon (1993), Estimation and Inference in Econometrics (general, a bit advanced)
7. Ruud (2000), Introduction to Classical Econometric Theory (general, consistent projection approach, careful)
8. Davidson (2000), Econometric Theory (econometrics/time series, LSE approach)
9. Mittelhammer, Judge, and Miller (2000), Econometric Foundations (general, advanced)
10. Patterson (2000), An Introduction to Applied Econometrics (econometrics/time series, LSE approach with applications)
11. Judge et al (1985), Theory and Practice of Econometrics (general, a bit old)
12. Hamilton (1994), Time Series Analysis
13. Spanos (1986), Statistical Foundations of Econometric Modelling, Cambridge University Press (general econometrics, LSE approach)
14. Harvey (1981), Time Series Models, Philip Allan
15. Harvey (1989), Forecasting, Structural Time Series... (structural time series, Kalman filter).
16. Lütkepohl (1993), Introduction to Multiple Time Series Analysis (time series, VAR models)
17. Priestley (1981), Spectral Analysis and Time Series (advanced time series)
18. Amemiya (1985), Advanced Econometrics, (asymptotic theory, non-linear econometrics)
19. Silverman (1986), Density Estimation for Statistics and Data Analysis (density estimation).
20. Härdle (1990), Applied Nonparametric Regression

## B A CLT in Action

This is an example of how we can calculate the limiting distribution of a sample average.
Remark 1 If $\sqrt{T}(\bar{x}-\mu) / \sigma \sim N(0,1)$ then $\bar{x} \sim N\left(\mu, \sigma^{2} / T\right)$.
Example 2 (Distribution of $\Sigma_{t=1}^{T}\left(z_{t}-1\right) / T$ and $\sqrt{T} \Sigma_{t=1}^{T}\left(z_{t}-1\right) / T$ when $z_{t} \sim \chi^{2}(1)$.) When $z_{t}$ is iid $\chi^{2}(1)$, then $\Sigma_{t=1}^{T} z_{t}$ is distributed as a $\chi^{2}(T)$ variable with pdf $f_{T}()$. We now construct a new variable by transforming $\Sigma_{t=1}^{T} z_{t}$ as to a sample mean around one (the mean of $z_{t}$ )

$$
\bar{z}_{1}=\Sigma_{t=1}^{T} z_{t} / T-1=\Sigma_{t=1}^{T}\left(z_{t}-1\right) / T
$$

Clearly, the inverse function is $\Sigma_{t=1}^{T} z_{t}=T \bar{z}_{1}+T$, so by the "change of variable" rule we get the pdf of $\bar{z}_{1}$ as

$$
g\left(\bar{z}_{1}\right)=f_{T}\left(T \bar{z}_{1}+T\right) T
$$

Example 3 Continuing the previous example, we now consider the random variable

$$
\bar{z}_{2}=\sqrt{T} \bar{z}_{1},
$$

with inverse function $\bar{z}_{1}=\bar{z}_{2} / \sqrt{T}$. By applying the "change of variable" rule again, we get the pdf of $\bar{z}_{2}$ as

$$
h\left(\bar{z}_{2}\right)=g\left(\bar{z}_{2} / \sqrt{T}\right) / \sqrt{T}=f_{T}\left(\sqrt{T} \bar{z}_{2}+T\right) \sqrt{T}
$$

Example 4 When $z_{t}$ is iid $\chi^{2}(1)$, then $\Sigma_{t=1}^{T} z_{t}$ is $\chi^{2}(T)$, which we denote $f\left(\Sigma_{t=1}^{T} z_{t}\right)$. We now construct two new variables by transforming $\Sigma_{t=1}^{T} z_{t}$

$$
\begin{aligned}
& \bar{z}_{1}=\Sigma_{t=1}^{T} z_{t} / T-1=\Sigma_{t=1}^{T}\left(z_{t}-1\right) / T, \text { and } \\
& \bar{z}_{2}=\sqrt{T} \bar{z}_{1}
\end{aligned}
$$

Example 5 We transform this distribution by first subtracting one from $z_{t}$ (to remove the mean) and then by dividing by $T$ or $\sqrt{T}$. This gives the distributions of the sample mean and scaled sample mean, $\bar{z}_{2}=\sqrt{T} \bar{z}_{1}$ as

$$
\begin{aligned}
& f\left(\bar{z}_{1}\right)=\frac{1}{2^{T / 2} \Gamma(T / 2)} y^{T / 2-1} \exp (-y / 2) \text { with } y=T \bar{z}_{1}+T, \text { and } \\
& f\left(\bar{z}_{2}\right)=\frac{1}{2^{T / 2} \Gamma(T / 2)} y^{T / 2-1} \exp (-y / 2) \text { with } y=\sqrt{T} \bar{z}_{1}+T
\end{aligned}
$$



Figure B.1: Sampling distributions. This figure shows the distribution of the sample mean and of $\sqrt{T}$ times the sample mean of the random variable $z_{t}-1$ where $z_{t} \sim \chi^{2}$ (1).

These distributions are shown in Figure B.1. It is clear that $f\left(\bar{z}_{1}\right)$ converges to a spike at zero as the sample size increases, while $f\left(\bar{z}_{2}\right)$ converges to a (non-trivial) normal distribution.

Example 6 (Distribution of $\Sigma_{t=1}^{T}\left(z_{t}-1\right) / T$ and $\sqrt{T} \Sigma_{t=1}^{T}\left(z_{t}-1\right) / T$ when $z_{t} \sim \chi^{2}(1)$.) When $z_{t}$ is iid $\chi^{2}(1)$, then $\Sigma_{t=1}^{T} z_{t}$ is $\chi^{2}(T)$, that is, has the probability density function

$$
f\left(\Sigma_{t=1}^{T} z_{t}\right)=\frac{1}{2^{T / 2} \Gamma(T / 2)}\left(\Sigma_{t=1}^{T} z_{t}\right)^{T / 2-1} \exp \left(-\Sigma_{t=1}^{T} z_{t} / 2\right)
$$

We transform this distribution by first subtracting one from $z_{t}$ (to remove the mean) and then by dividing by $T$ or $\sqrt{T}$. This gives the distributions of the sample mean, $\bar{z}_{1}=$ $\Sigma_{t=1}^{T}\left(z_{t}-1\right) / T$, and scaled sample mean, $\bar{z}_{2}=\sqrt{T} \bar{z}_{1}$ as

$$
\begin{aligned}
& f\left(\bar{z}_{1}\right)=\frac{1}{2^{T / 2} \Gamma(T / 2)} y^{T / 2-1} \exp (-y / 2) \text { with } y=T \bar{z}_{1}+T, \text { and } \\
& f\left(\bar{z}_{2}\right)=\frac{1}{2^{T / 2} \Gamma(T / 2)} y^{T / 2-1} \exp (-y / 2) \text { with } y=\sqrt{T} \bar{z}_{1}+T
\end{aligned}
$$

These distributions are shown in Figure B.1. It is clear that $f\left(\bar{z}_{1}\right)$ converges to a spike at zero as the sample size increases, while $f\left(\bar{z}_{2}\right)$ converges to a (non-trivial) normal distribution.

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Silverman, B. W., 1986, Density Estimation for Statistics and Data Analysis, Chapman and Hall, London.

Verbeek, M., 2000, A Guide to Modern Econometrics, Wiley, Chichester.

## 2 Univariate Time Series Analysis

Reference: Greene (2000) 13.1-3 and 18.1-3
Additional references: Hayashi (2000) 6.2-4; Verbeek (2000) 8-9; Hamilton (1994); Johnston and DiNardo (1997) 7; and Pindyck and Rubinfeld (1997) 16-18

### 2.1 Theoretical Background to Time Series Processes

Suppose we have a sample of $T$ observations of a random variable

$$
\left\{y_{t}^{i}\right\}_{t=1}^{T}=\left\{y_{1}^{i}, y_{2}^{i}, \ldots, y_{T}^{i}\right\},
$$

where subscripts indicate time periods. The superscripts indicate that this sample is from planet (realization) $i$. We could imagine a continuum of parallel planets where the same time series process has generated different samples with $T$ different numbers (different realizations).

Consider period $t$. The distribution of $y_{t}$ across the (infinite number of) planets has some density function, $f_{t}\left(y_{t}\right)$. The mean of this distribution

$$
\begin{equation*}
\mathrm{E} y_{t}=\int_{-\infty}^{\infty} y_{t} f_{t}\left(y_{t}\right) d y_{t} \tag{2.1}
\end{equation*}
$$

is the expected value of the value in period $t$, also called the unconditional mean of $y_{t}$. Note that $\mathrm{E} y_{t}$ could be different from $\mathrm{E} y_{t+s}$. The unconditional variance is defined similarly.

Now consider periods $t$ and $t-s$ jointly. On planet $i$ we have the pair $\left\{y_{t-s}^{i}, y_{t}^{i}\right\}$. The bivariate distribution of these pairs, across the planets, has some density function $g_{t-s, t}\left(y_{t-s}, y_{t}\right) .{ }^{1}$ Calculate the covariance between $y_{t-s}$ and $y_{t}$ as usual

$$
\begin{align*}
\operatorname{Cov}\left(y_{t-s}, y_{t}\right) & =\int_{-\infty}^{\infty} \int_{-\infty}^{\infty}\left(y_{t-s}-\mathrm{E} y_{t-s}\right)\left(y_{t}-\mathrm{E} y_{t}\right) g_{t-s, t}\left(y_{t-s}, y_{t}\right) d y_{t} d y_{t-s}  \tag{2.2}\\
& =\mathrm{E}\left(y_{t-s}-\mathrm{E} y_{t-s}\right)\left(y_{t}-\mathrm{E} y_{t}\right) \tag{2.3}
\end{align*}
$$

[^0]This is the $s^{\text {th }}$ autocovariance of $y_{t}$. (Of course, $s=0$ or $s<0$ are allowed.)
A stochastic process is covariance stationary if

$$
\begin{align*}
& \mathrm{E} y_{t}=\mu \text { is independent of } t,  \tag{2.4}\\
& \operatorname{Cov}\left(y_{t-s}, y_{t}\right)=\gamma_{s} \text { depends only on } s, \text { and }  \tag{2.5}\\
& \text { both } \mu \text { and } \gamma_{s} \text { are finite. } \tag{2.6}
\end{align*}
$$

Most of these notes are about covariance stationary processes, but Section 2.7 is about non-stationary processes.

Humanity has so far only discovered one planet with coin flipping; any attempt to estimate the moments of a time series process must therefore be based on the realization of the stochastic process from planet earth only. This is meaningful only if the process is ergodic for the moment you want to estimate. A covariance stationary process is said to be ergodic for the mean if

$$
\begin{equation*}
\operatorname{plim} \frac{1}{T} \sum_{t=1}^{T} y_{t}=\mathrm{E} y_{t} \tag{2.7}
\end{equation*}
$$

so the sample mean converges in probability to the unconditional mean. A sufficient condition for ergodicity for the mean is

$$
\begin{equation*}
\sum_{s=0}^{\infty}\left|\operatorname{Cov}\left(y_{t-s}, y_{t}\right)\right|<\infty \tag{2.8}
\end{equation*}
$$

This means that the link between the values in $t$ and $t-s$ goes to zero sufficiently fast as $s$ increases (you may think of this as getting independent observations before we reach the limit). If $y_{t}$ is normally distributed, then (2.8) is also sufficient for the process to be ergodic for all moments, not just the mean. Figure 2.1 illustrates how a longer and longer sample (of one realization of the same time series process) gets closer and closer to the unconditional distribution as the sample gets longer.

### 2.2 Estimation of Autocovariances

Let $y_{t}$ be a vector of a covariance stationary and ergodic. The $s$ th covariance matrix is

$$
\begin{equation*}
R(s)=\mathrm{E}\left(y_{t}-\mathrm{E} y_{t}\right)\left(y_{t-s}-\mathrm{E} y_{t-s}\right)^{\prime} . \tag{2.9}
\end{equation*}
$$




Mean and Std over longer and longer samples


Figure 2.1: Sample of one realization of $y_{t}=0.85 y_{t-1}+\varepsilon_{t}$ with $y_{0}=4$ and $\operatorname{Std}\left(\varepsilon_{t}\right)=1$.

Note that $R(s)$ does not have to be symmetric unless $s=0$. However, note that $R(s)=$ $R(-s)^{\prime}$. This follows from noting that

$$
\begin{align*}
R(-s) & =\mathrm{E}\left(y_{t}-\mathrm{E} y_{t}\right)\left(y_{t+s}-\mathrm{E} y_{t+s}\right)^{\prime} \\
& =\mathrm{E}\left(y_{t-s}-\mathrm{E} y_{t-s}\right)\left(y_{t}-\mathrm{E} y_{t}\right)^{\prime}, \tag{2.10a}
\end{align*}
$$

where we have simply changed time subscripts and exploited the fact that $y_{t}$ is covariance stationary. Transpose to get

$$
\begin{equation*}
R(-s)^{\prime}=\mathrm{E}\left(y_{t}-\mathrm{E} y_{t}\right)\left(y_{t-s}-\mathrm{E} y_{t-s}\right)^{\prime}, \tag{2.11}
\end{equation*}
$$

which is the same as in (2.9). If $y_{t}$ is a scalar, then $R(s)=R(-s)$, which shows that autocovariances are symmetric around $s=0$.

Example 1 (Bivariate case.) Let $y_{t}=\left[x_{t}, z_{t}\right]^{\prime}$ with $E x_{t}=E z_{t}=0$. Then

$$
\begin{aligned}
\hat{R}(s) & =E\left[\begin{array}{l}
x_{t} \\
z_{t}
\end{array}\right]\left[\begin{array}{ll}
x_{t-s} & z_{t-s}
\end{array}\right] \\
& =\left[\begin{array}{ll}
\operatorname{Cov}\left(x_{t}, x_{t-s}\right) & \operatorname{Cov}\left(x_{t}, z_{t-s}\right) \\
\operatorname{Cov}\left(z_{t}, x_{t-s}\right) & \operatorname{Cov}\left(z_{t}, x_{t-s}\right)
\end{array}\right] .
\end{aligned}
$$

Note that $R(-s)$ is

$$
\begin{aligned}
R(-s) & =\left[\begin{array}{ll}
\operatorname{Cov}\left(x_{t}, x_{t+s}\right) & \operatorname{Cov}\left(x_{t}, z_{t+s}\right) \\
\operatorname{Cov}\left(z_{t}, x_{t+s}\right) & \operatorname{Cov}\left(z_{t}, x_{t+s}\right)
\end{array}\right] \\
& =\left[\begin{array}{ll}
\operatorname{Cov}\left(x_{t-s}, x_{t}\right) & \operatorname{Cov}\left(x_{t-s}, z_{t}\right) \\
\operatorname{Cov}\left(z_{t-s}, x_{t}\right) & \operatorname{Cov}\left(z_{t-s}, x_{t}\right)
\end{array}\right],
\end{aligned}
$$

which is indeed the transpose of $R(s)$.
The autocovariances of the (vector) $y_{t}$ process can be estimated as

$$
\begin{align*}
\hat{R}(s) & =\frac{1}{T} \sum_{t=1+s}^{T}\left(y_{t}-\bar{y}\right)\left(y_{t-s}-\bar{y}\right)^{\prime}  \tag{2.12}\\
\text { with } \bar{y} & =\frac{1}{T} \sum_{t=1}^{T} y_{t} \tag{2.13}
\end{align*}
$$

(We typically divide by $T$ in even if we have only $T-s$ full observations to estimate $R(s)$ from.)

Autocorrelations are then estimated by dividing the diagonal elements in $\hat{R}(s)$ by the diagonal elements in $\hat{R}(0)$

$$
\begin{equation*}
\hat{\rho}(s)=\operatorname{diag} \hat{R}(s) / \operatorname{diag} \hat{R}(0) \text { (element by element }) \tag{2.14}
\end{equation*}
$$

### 2.3 White Noise

A white noise time process has

$$
\begin{align*}
\mathrm{E} \varepsilon_{t} & =0 \\
\operatorname{Var}\left(\varepsilon_{t}\right) & =\sigma^{2}, \text { and } \\
\operatorname{Cov}\left(\varepsilon_{t-s}, \varepsilon_{t}\right) & =0 \text { if } s \neq 0 . \tag{2.15}
\end{align*}
$$

If, in addition, $\varepsilon_{t}$ is normally distributed, then it is said to be Gaussian white noise. The conditions in (2.4)-(2.6) are satisfied so this process is covariance stationary. Moreover, (2.8) is also satisfied, so the process is ergodic for the mean (and all moments if $\varepsilon_{t}$ is normally distributed).

### 2.4 Moving Average

A $q^{\text {th }}$-order moving average process is

$$
\begin{equation*}
y_{t}=\varepsilon_{t}+\theta_{1} \varepsilon_{t-1}+\ldots+\theta_{q} \varepsilon_{t-q}, \tag{2.16}
\end{equation*}
$$

where the innovation $\varepsilon_{t}$ is white noise (usually Gaussian). We could also allow both $y_{t}$ and $\varepsilon_{t}$ to be vectors; such a process it called a vector MA (VMA).

We have $\mathrm{E} y_{t}=0$ and

$$
\begin{align*}
\operatorname{Var}\left(y_{t}\right) & =\mathrm{E}\left(\varepsilon_{t}+\theta_{1} \varepsilon_{t-1}+\ldots+\theta_{q} \varepsilon_{t-q}\right)\left(\varepsilon_{t}+\theta_{1} \varepsilon_{t-1}+\ldots+\theta_{q} \varepsilon_{t-q}\right) \\
& =\sigma^{2}\left(1+\theta_{1}^{2}+\ldots+\theta_{q}^{2}\right) . \tag{2.17}
\end{align*}
$$

Autocovariances are calculated similarly, and it should be noted that autocovariances of order $q+1$ and higher are always zero for an MA $(q)$ process.

Example 2 The mean of an $M A(1), y_{t}=\varepsilon_{t}+\theta_{1} \varepsilon_{t-1}$, is zero since the mean of $\varepsilon_{t}$ (and $\left.\varepsilon_{t-1}\right)$ is zero. The first three autocovariance are

$$
\begin{align*}
\operatorname{Var}\left(y_{t}\right) & =E\left(\varepsilon_{t}+\theta_{1} \varepsilon_{t-1}\right)\left(\varepsilon_{t}+\theta_{1} \varepsilon_{t-1}\right)=\sigma^{2}\left(1+\theta_{1}^{2}\right) \\
\operatorname{Cov}\left(y_{t-1}, y_{t}\right) & =E\left(\varepsilon_{t-1}+\theta_{1} \varepsilon_{t-2}\right)\left(\varepsilon_{t}+\theta_{1} \varepsilon_{t-1}\right)=\sigma^{2} \theta_{1} \\
\operatorname{Cov}\left(y_{t-2}, y_{t}\right) & =E\left(\varepsilon_{t-2}+\theta_{1} \varepsilon_{t-3}\right)\left(\varepsilon_{t}+\theta_{1} \varepsilon_{t-1}\right)=0, \tag{2.18}
\end{align*}
$$

and $\operatorname{Cov}\left(y_{t-s}, y_{t}\right)=0$ for $|s| \geq 2$. Since both the mean and the covariances are finite and constant across $t$, the $M A(1)$ is covariance stationary. Since the absolute value of the covariances sum to a finite number, the $M A(1)$ is also ergodic for the mean. The first autocorrelation of an MA(1) is

$$
\operatorname{Corr}\left(y_{t-1}, y_{t}\right)=\frac{\theta_{1}}{1+\theta_{1}^{2}}
$$

Since the white noise process is covariance stationary, and since an MA $(q)$ with $m<$ $\infty$ is a finite order linear function of $\varepsilon_{t}$, it must be the case that the MA $(q)$ is covariance stationary. It is ergodic for the mean since $\operatorname{Cov}\left(y_{t-s}, y_{t}\right)=0$ for $s>q$, so (2.8) is satisfied. As usual, Gaussian innovations are then sufficient for the MA $(q)$ to be ergodic for all moments.

The effect of $\varepsilon_{t}$ on $y_{t}, y_{t+1}, \ldots$, that is, the impulse response function, is the same as the MA coefficients

$$
\begin{equation*}
\frac{\partial y_{t}}{\partial \varepsilon_{t}}=1, \frac{\partial y_{t+1}}{\partial \varepsilon_{t}}=\theta_{1}, \ldots, \frac{\partial y_{t+q}}{\partial \varepsilon_{t}}=\theta_{q}, \text { and } \frac{\partial y_{t+q+k}}{\partial \varepsilon_{t}}=0 \text { for } k>0 . \tag{2.19}
\end{equation*}
$$

This is easily seen from applying (2.16)

$$
\begin{aligned}
& y_{t}=\underline{\varepsilon_{t}}+\theta_{1} \varepsilon_{t-1}+\ldots+\theta_{q} \varepsilon_{t-q} \\
& y_{t+1}=\varepsilon_{t+1}+\theta_{1} \underline{\varepsilon_{t}}+\ldots+\theta_{q} \varepsilon_{t-q+1} \\
& \vdots \\
& y_{t+q}=\varepsilon_{t+q}+\theta_{1} \varepsilon_{t-1+q}+\ldots+\theta_{q} \varepsilon_{t} \\
& y_{t+q+1}=\varepsilon_{t+q+1}+\theta_{1} \varepsilon_{t+q}+\ldots+\theta_{q} \varepsilon_{t+1} .
\end{aligned}
$$

The expected value of $y_{t}$, conditional on $\left\{\varepsilon_{w}\right\}_{w=-\infty}^{t-s}$ is

$$
\begin{align*}
\mathrm{E}_{t-s} y_{t} & =\mathrm{E}_{t-s}\left(\varepsilon_{t}+\theta_{1} \varepsilon_{t-1}+\ldots+\theta_{s} \varepsilon_{t-s}+\ldots+\theta_{q} \varepsilon_{t-q}\right) \\
& =\theta_{s} \varepsilon_{t-s}+\ldots+\theta_{q} \varepsilon_{t-q}, \tag{2.20}
\end{align*}
$$

since $\mathrm{E}_{t-s} \varepsilon_{t-(s-1)}=\ldots=\mathrm{E}_{t-s} \varepsilon_{t}=0$.
Example 3 (Forecasting an MA(1).) Suppose the process is

$$
y_{t}=\varepsilon_{t}+\theta_{1} \varepsilon_{t-1}, \text { with } \operatorname{Var}\left(\varepsilon_{t}\right)=\sigma^{2} .
$$

The forecasts made in $t=2$ then have the follow expressions-with an example using $\theta_{1}=2, \varepsilon_{1}=3 / 4$ and $\varepsilon_{2}=1 / 2$ in the second column

$$
\begin{array}{lr}
\text { General } & \frac{\text { Example }}{\text { yen }} \\
y_{2} & =1 / 2+2 \times 3 / 4=2 \\
E_{2} y_{3}=E_{2}\left(\varepsilon_{3}+\theta_{1} \varepsilon_{2}\right)=\theta_{1} \varepsilon_{2} & =2 \times 1 / 2=1 \\
E_{2} y_{4}=E_{2}\left(\varepsilon_{4}+\theta_{1} \varepsilon_{3}\right)=0 & =0
\end{array}
$$

Example 4 (MA(1) and conditional variances.) From Example 3, the forecasting variances are—with the numerical example continued assuming that $\sigma^{2}=1$

| General | $\frac{\text { Example }}{}$ |
| :--- | ---: |
| $\operatorname{Var}\left(y_{2}-E_{2} y_{2}\right)=0$ | $=0$ |
| $\operatorname{Var}\left(y_{3}-E_{2} y_{3}\right)=\operatorname{Var}\left(\varepsilon_{3}+\theta_{1} \varepsilon_{2}-\theta_{1} \varepsilon_{2}\right)=\sigma^{2}$ | $=1$ |
| $\operatorname{Var}\left(y_{4}-E_{2} y_{4}\right)=\operatorname{Var}\left(\varepsilon_{4}+\theta_{1} \varepsilon_{3}\right)=\sigma^{2}+\theta_{1}^{2} \sigma^{2}$ | $=5$ |

If the innovations are iid Gaussian, then the distribution of the $s$-period forecast error

$$
y_{t}-\mathrm{E}_{t-s} y_{t}=\varepsilon_{t}+\theta_{1} \varepsilon_{t-1}+\ldots+\theta_{s-1} \varepsilon_{t-(s-1)}
$$

is

$$
\begin{equation*}
\left(y_{t}-\mathrm{E}_{t-s} y_{t}\right) \sim N\left[0, \sigma^{2}\left(1+\theta_{1}^{2}+\ldots+\theta_{s-1}^{2}\right)\right] \tag{2.21}
\end{equation*}
$$

since $\varepsilon_{t}, \varepsilon_{t-1}, \ldots, \varepsilon_{t-(s-1)}$ are independent Gaussian random variables. This implies that the conditional distribution of $y_{t}$, conditional on $\left\{\varepsilon_{w}\right\}_{w}^{s}$ $\qquad$

$$
\begin{align*}
y_{t} \mid\left\{\varepsilon_{t-s}, \varepsilon_{t-s-1}, \ldots\right\} & \sim N\left[\mathrm{E}_{t-s} y_{t}, \operatorname{Var}\left(y_{t}-\mathrm{E}_{t-s} y_{t}\right)\right]  \tag{2.22}\\
& \sim N\left[\theta_{s} \varepsilon_{t-s}+\ldots+\theta_{q} \varepsilon_{t-q}, \sigma^{2}\left(1+\theta_{1}^{2}+\ldots+\theta_{s-1}^{2}\right)\right] . \tag{2.23}
\end{align*}
$$

The conditional mean is the point forecast and the variance is the variance of the forecast error. Note that if $s>q$, then the conditional distribution coincides with the unconditional distribution since $\varepsilon_{t-s}$ for $s>q$ is of no help in forecasting $y_{t}$.

Example 5 (MA(1) and convergence from conditional to unconditional distribution.) From examples 3 and 4 we see that the conditional distributions change according to (where
$\Omega_{2}$ indicates the information set in $t=2$ )

| General | $\left.=\begin{array}{l}\text { Example } \\ \hline y_{2} \mid \Omega_{2} \sim N\left(y_{2}, 0\right) \\ y_{3} \mid \Omega_{2} \sim N\left(E_{2} y_{3}, \operatorname{Var}\left(y_{3}-E_{2} y_{3}\right)\right) \\ y_{4} \mid \Omega_{2} \sim N\left(E_{2} y_{4}, \operatorname{Var}\left(y_{4}-E_{2} y_{4}\right)\right) \\ \hline\end{array}\right) N(0,5)$ |
| :--- | :--- |

Note that the distribution of $y_{4} \mid \Omega_{2}$ coincides with the asymptotic distribution.

Estimation of MA processes is typically done by setting up the likelihood function and then using some numerical method to maximize it.

### 2.5 Autoregression

A $p^{t h}$-order autoregressive process is

$$
\begin{equation*}
y_{t}=a_{1} y_{t-1}+a_{2} y_{t-2}+\ldots+a_{p} y_{t-p}+\varepsilon_{t} . \tag{2.24}
\end{equation*}
$$

$\operatorname{A~} \operatorname{VAR}(p)$ is just like the $\operatorname{AR}(p)$ in (2.24), but where $y_{t}$ is interpreted as a vector and $a_{i}$ as a matrix.

Example 6 (VAR(1) model.) A VAR(1) model is of the following form

$$
\left[\begin{array}{l}
y_{1 t} \\
y_{2 t}
\end{array}\right]=\left[\begin{array}{ll}
a_{11} & a_{12} \\
a_{21} & a_{22}
\end{array}\right]\left[\begin{array}{l}
y_{1 t-1} \\
y_{2 t-1}
\end{array}\right]+\left[\begin{array}{l}
\varepsilon_{1 t} \\
\varepsilon_{2 t}
\end{array}\right]
$$

All stationary $\operatorname{AR}(p)$ processes can be written on $\mathrm{MA}(\infty)$ form by repeated substitution. To do so we rewrite the $\operatorname{AR}(p)$ as a first order vector autoregression, $\operatorname{VAR}(1)$. For instance, an $\operatorname{AR}(2) x_{t}=a_{1} x_{t-1}+a_{2} x_{t-2}+\varepsilon_{t}$ can be written as

$$
\begin{align*}
{\left[\begin{array}{l}
x_{t} \\
x_{t-1}
\end{array}\right] } & =\left[\begin{array}{ll}
a_{1} & a_{2} \\
1 & 0
\end{array}\right]\left[\begin{array}{l}
x_{t-1} \\
x_{t-2}
\end{array}\right]+\left[\begin{array}{l}
\varepsilon_{t} \\
0
\end{array}\right], \text { or }  \tag{2.25}\\
y_{t} & =A y_{t-1}+\varepsilon_{t}, \tag{2.26}
\end{align*}
$$

where $y_{t}$ is an $2 \times 1$ vector and $A$ a $4 \times 4$ matrix. This works also if $x_{t}$ and $\varepsilon_{t}$ are vectors and. In this case, we interpret $a_{i}$ as matrices and 1 as an identity matrix.

Iterate backwards on (2.26)

$$
\begin{align*}
y_{t} & =A\left(A y_{t-2}+\varepsilon_{t-1}\right)+\varepsilon_{t} \\
& =A^{2} y_{t-2}+A \varepsilon_{t-1}+\varepsilon_{t} \\
\vdots &  \tag{2.27}\\
& =A^{K+1} y_{t-K-1}+\sum_{s=0}^{K} A^{s} \varepsilon_{t-s} .
\end{align*}
$$

Remark 7 (Spectral decomposition.) The $n$ eigenvalues $\left(\lambda_{i}\right)$ and associated eigenvectors $\left(z_{i}\right)$ of the $n \times n$ matrix A satisfy

$$
\left(A-\lambda_{i} I_{n}\right) z_{i}=\mathbf{0}_{n \times 1} .
$$

If the eigenvectors are linearly independent, then

$$
A=Z \Lambda Z^{-1} \text {, where } \Lambda=\left[\begin{array}{cccc}
\lambda_{1} & 0 & \cdots & 0 \\
0 & \lambda_{2} & \cdots & 0 \\
\vdots & \vdots & \cdots & \vdots \\
0 & 0 & \cdots & \lambda_{n}
\end{array}\right] \text { and } Z=\left[\begin{array}{llll}
z_{1} & z_{2} & \cdots & z_{n}
\end{array}\right]
$$

Note that we therefore get

$$
A^{2}=A A=Z \Lambda Z^{-1} Z \Lambda Z^{-1}=Z \Lambda \Lambda Z^{-1}=Z \Lambda^{2} Z^{-1} \Rightarrow A^{q}=Z \Lambda^{q} Z^{-1}
$$

Remark 8 (Modulus of complex number.) If $\lambda=a+$ bi, where $i=\sqrt{-1}$, then $|\lambda|=$ $|a+b i|=\sqrt{a^{2}+b^{2}}$.

Take the limit of (2.27) as $K \rightarrow \infty$. If $\lim _{K \rightarrow \infty} A^{K+1} y_{t-K-1}=0$, then we have a moving average representation of $y_{t}$ where the influence of the starting values vanishes asymptotically

$$
\begin{equation*}
y_{t}=\sum_{s=0}^{\infty} A^{s} \varepsilon_{t-s} . \tag{2.28}
\end{equation*}
$$

We note from the spectral decompositions that $A^{K+1}=Z \Lambda^{K+1} Z^{-1}$, where $Z$ is the matrix of eigenvectors and $\Lambda$ a diagonal matrix with eigenvalues. Clearly, $\lim _{K \rightarrow \infty} A^{K+1} y_{t-K-1}=$ 0 is satisfied if the eigenvalues of $A$ are all less than one in modulus and $y_{t-K-1}$ does not grow without a bound.


Figure 2.2: Conditional moments and distributions for different forecast horizons for the $\operatorname{AR}(1)$ process $y_{t}=0.85 y_{t-1}+\varepsilon_{t}$ with $y_{0}=4$ and $\operatorname{Std}\left(\varepsilon_{t}\right)=1$.

Example $9\left(A R(1)\right.$.) For the univariate $A R(1) y_{t}=a y_{t-1}+\varepsilon_{t}$, the characteristic equation is $(a-\lambda) z=0$, which is only satisfied if the eigenvalue is $\lambda=a$. The $A R(1)$ is therefore stable (and stationarity) if $-1<a<1$. This can also be seen directly by noting that $a^{K+1} y_{t-K-1}$ declines to zero if $0<a<1$ as $K$ increases.

Similarly, most finite order MA processes can be written ("inverted") as $\operatorname{AR}(\infty)$. It is therefore common to approximate MA processes with AR processes, especially since the latter are much easier to estimate.

Example 10 (Variance of $A R(1)$.) From the $M A$-representation $y_{t}=\sum_{s=0}^{\infty} a^{s} \varepsilon_{t-s}$ and the fact that $\varepsilon_{t}$ is white noise we get $\operatorname{Var}\left(y_{t}\right)=\sigma^{2} \sum_{s=0}^{\infty} a^{2 s}=\sigma^{2} /\left(1-a^{2}\right)$. Note that this is minimized at $a=0$. The autocorrelations are obviously $a^{|s|}$. The covariance matrix of $\left\{y_{t}\right\}_{t=1}^{T}$ is therefore (standard deviation $\times$ standard deviation $\times$ autocorrelation)

$$
\frac{\sigma^{2}}{1-a^{2}}\left[\begin{array}{lllll}
1 & a & a^{2} & \cdots & a^{T-1} \\
a & 1 & a & \cdots & a^{T-2} \\
a^{2} & a & 1 & \cdots & a^{T-3} \\
\vdots & & & \ddots & \\
a^{T-1} & a^{T-2} & a^{T-3} & \cdots & 1
\end{array}\right]
$$

Example 11 (Covariance stationarity of an $A R(1)$ with $|a|<1$.) From the MA-representation $y_{t}=\sum_{s=0}^{\infty} a^{s} \varepsilon_{t-s}$, the expected value of $y_{t}$ is zero, since $E \varepsilon_{t-s}=0$. We know that $\operatorname{Cov}\left(y_{t}, y_{t-s}\right)=a^{|s|} \sigma^{2} /\left(1-a^{2}\right)$ which is constant and finite.

Example 12 (Ergodicity of a stationary $A R(1)$.) We know that $\operatorname{Cov}\left(y_{t}, y_{t-s}\right)=a^{|s|} \sigma^{2} /\left(1-a^{2}\right)$, so the absolute value is

$$
\left|\operatorname{Cov}\left(y_{t}, y_{t-s}\right)\right|=|a|^{|s|} \sigma^{2} /\left(1-a^{2}\right)
$$

Using this in (2.8) gives

$$
\begin{aligned}
\sum_{s=0}^{\infty}\left|\operatorname{Cov}\left(y_{t-s}, y_{t}\right)\right| & =\frac{\sigma^{2}}{1-a^{2}} \sum_{s=0}^{\infty}|a|^{s} \\
& =\frac{\sigma^{2}}{1-a^{2}} \frac{1}{1-|a|}(\text { since }|a|<1)
\end{aligned}
$$

which is finite. The $A R(1)$ is ergodic if $|a|<1$.
Example 13 (Conditional distribution of $A R(1)$.) For the $A R(1) y_{t}=a y_{t-1}+\varepsilon_{t}$ with $\varepsilon_{t} \sim N\left(0, \sigma^{2}\right)$, we get

$$
\begin{aligned}
E_{t} y_{t+s} & =a^{s} y_{t} \\
\operatorname{Var}\left(y_{t+s}-E_{t} y_{t+s}\right) & =\left(1+a^{2}+a^{4}+\ldots+a^{2(s-1)}\right) \sigma^{2} \\
& =\frac{a^{2 s}-1}{a^{2}-1} \sigma^{2}
\end{aligned}
$$

The distribution of $y_{t+s}$ conditional on $y_{t}$ is normal with these parameters. See Figure 2.2 for an example.

### 2.5.1 Estimation of an AR(1) Process

Suppose we have sample $\left\{y_{t}\right\}_{t=0}^{T}$ of a process which we know is an $\operatorname{AR}(p), y_{t}=a y_{t-1}+$ $\varepsilon_{t}$, with normally distributed innovations with unknown variance $\sigma^{2}$.

The pdf of $y_{1}$ conditional on $y_{0}$ is

$$
\begin{equation*}
\operatorname{pdf}\left(y_{1} \mid y_{0}\right)=\frac{1}{\sqrt{2 \pi \sigma^{2}}} \exp \left(-\frac{\left(y_{1}-a y_{0}\right)^{2}}{2 \sigma^{2}}\right) \tag{2.29}
\end{equation*}
$$

and the pdf of $y_{2}$ conditional on $y_{1}$ and $y_{0}$ is

$$
\begin{equation*}
\operatorname{pdf}\left(y_{2} \mid\left\{y_{1}, y_{0}\right\}\right)=\frac{1}{\sqrt{2 \pi \sigma^{2}}} \exp \left(-\frac{\left(y_{2}-a y_{1}\right)^{2}}{2 \sigma^{2}}\right) . \tag{2.30}
\end{equation*}
$$

### 2.5.2 Lag Operators*

A common and convenient way of dealing with leads and lags is the lag operator, L. It is such that

$$
\mathrm{L}^{s} y_{t}=y_{t-s} \text { for all (integer) } s
$$

For instance, the ARMA(2,1) model

$$
\begin{equation*}
y_{t}-a_{1} y_{t-1}-a_{2} y_{t-2}=\varepsilon_{t}+\theta_{1} \varepsilon_{t-1} \tag{2.35}
\end{equation*}
$$

can be written as

$$
\begin{equation*}
\left(1-a_{1} \mathrm{~L}-a_{2} \mathrm{~L}^{2}\right) y_{t}=\left(1+\theta_{1} \mathrm{~L}\right) \varepsilon_{t} \tag{2.36}
\end{equation*}
$$

which is usually denoted

$$
\begin{equation*}
a(\mathrm{~L}) y_{t}=\theta(\mathrm{L}) \varepsilon_{t} \tag{2.37}
\end{equation*}
$$

### 2.5.3 Properties of LS Estimates of an $\operatorname{AR}(p)$ Process*

## Reference: Hamilton (1994) 8.2

The LS estimates are typically biased, but consistent and asymptotically normally distributed, provided the AR is stationary.

As usual the LS estimate is

$$
\left.\begin{array}{rl}
\hat{\beta}_{L S}-\beta & =\left[\frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right.
\end{array}\right]^{-1} \frac{1}{T} \sum_{t=1}^{T} x_{t} \varepsilon_{t}, \text { where } ~ 子 x_{t}=\left[\begin{array}{llll}
y_{t-1} & y_{t-2} & \cdots & y_{t-p} \tag{2.38}
\end{array}\right] .
$$

The first term in (2.38) is the inverse of the sample estimate of covariance matrix of $x_{t}$ (since $\mathrm{E} y_{t}=0$ ), which converges in probability to $\Sigma_{x x}^{-1}\left(y_{t}\right.$ is stationary and ergodic for all moments if $\varepsilon_{t}$ is Gaussian). The last term, $\frac{1}{T} \sum_{t=1}^{T} x_{t} \varepsilon_{t}$, is serially uncorrelated, so we can apply a CLT. Note that $\mathrm{E} x_{t} \varepsilon_{t} \varepsilon_{t}^{\prime} x_{t}^{\prime}=\mathrm{E} \varepsilon_{t} \varepsilon_{t}^{\prime} \mathrm{E} x_{t} x_{t}^{\prime}=\sigma^{2} \Sigma_{x x}$ since $u_{t}$ and $x_{t}$ are independent. We therefore have

$$
\begin{equation*}
\frac{1}{\sqrt{T}} \sum_{t=1}^{T} x_{t} \varepsilon_{t} \rightarrow^{d} N\left(0, \sigma^{2} \Sigma_{x x}\right) \tag{2.39}
\end{equation*}
$$



Figure 2.3: Impulse response function of ARMA(1,1)

### 2.7 Non-stationary Processes

### 2.7.1 Introduction

A trend-stationary process can be made stationary by subtracting a linear trend. The simplest example is

$$
\begin{equation*}
y_{t}=\mu+\beta t+\varepsilon_{t} \tag{2.46}
\end{equation*}
$$

## where $\varepsilon_{t}$ is white noise.

A unit root process can be made stationary only by taking a difference. The simplest example is the random walk with drift

$$
\begin{equation*}
y_{t}=\mu+y_{t-1}+\varepsilon_{t} \tag{2.47}
\end{equation*}
$$

where $\varepsilon_{t}$ is white noise. The name "unit root process" comes from the fact that the largest
eigenvalues of the canonical form (the $\operatorname{VAR}(1)$ form of the $\operatorname{AR}(p)$ ) is one. Such a process is said to be integrated of order one (often denoted $\mathrm{I}(1)$ ) and can be made stationary by taking first differences.

Example 14 (Non-stationary $A R(2)$.) The process $y_{t}=1.5 y_{t-1}-0.5 y_{t-2}+\varepsilon_{t}$ can be written

$$
\left[\begin{array}{l}
y_{t} \\
y_{t-1}
\end{array}\right]=\left[\begin{array}{ll}
1.5 & -0.5 \\
1 & 0
\end{array}\right]\left[\begin{array}{l}
y_{t-1} \\
y_{t-2}
\end{array}\right]+\left[\begin{array}{l}
\varepsilon_{t} \\
0
\end{array}\right]
$$

where the matrix has the eigenvalues 1 and 0.5 and is therefore non-stationary. Note that subtracting $y_{t-1}$ from both sides gives $y_{t}-y_{t-1}=0.5\left(y_{t-1}-y_{t-2}\right)+\varepsilon_{t}$, so the variable $x_{t}=y_{t}-y_{t-1}$ is stationary.

The distinguishing feature of unit root processes is that the effect of a shock never vanishes. This is most easily seen for the random walk. Substitute repeatedly in (2.47) to get

$$
\begin{align*}
y_{t} & =\mu+\left(\mu+y_{t-2}+\varepsilon_{t-1}\right)+\varepsilon_{t} \\
& \vdots  \tag{2.48}\\
& =t \mu+y_{0}+\sum_{s=1}^{t} \varepsilon_{s} .
\end{align*}
$$

The effect of $\varepsilon_{t}$ never dies out: a non-zero value of $\varepsilon_{t}$ gives a permanent shift of the level of $y_{t}$. This process is clearly non-stationary. A consequence of the permanent effect of a shock is that the variance of the conditional distribution grows without bound as the forecasting horizon is extended. For instance, for the random walk with drift, (2.48), the distribution conditional on the information in $t=0$ is $N\left(y_{0}+t \mu, s \sigma^{2}\right)$ if the innovations are Gaussian. This means that the expected change is $t \mu$ and that the conditional variance grows linearly with the forecasting horizon. The unconditional variance is therefore infinite and the standard results on inference are not applicable.

In contrast, the conditional distributions from the trend stationary model, (2.46), is $N\left(s t, \sigma^{2}\right)$.

A process could have two unit roots (integrated of order 2: I(2)). In this case, we need to difference twice to make it stationary. Alternatively, a process can also be explosive, that is, have eigenvalues outside the unit circle. In this case, the impulse response function diverges.

Example 15 (Two unit roots.) Suppose $y_{t}$ in Example (14) is actually the first difference of some other series, $y_{t}=z_{t}-z_{t-1}$. We then have

$$
\begin{aligned}
z_{t}-z_{t-1} & =1.5\left(z_{t-1}-z_{t-2}\right)-0.5\left(z_{t-2}-z_{t-3}\right)+\varepsilon_{t} \\
z_{t} & =2.5 z_{t-1}-2 z_{t-2}+0.5 z_{t-3}+\varepsilon_{t},
\end{aligned}
$$

which is an $A R(3)$ with the following canonical form

$$
\left[\begin{array}{l}
z_{t} \\
z_{t-1} \\
z_{t-2}
\end{array}\right]=\left[\begin{array}{ccc}
2.5 & -2 & 0.5 \\
1 & 0 & 0 \\
0 & 1 & 0
\end{array}\right]\left[\begin{array}{c}
z_{t-1} \\
z_{t-2} \\
z_{t-3}
\end{array}\right]+\left[\begin{array}{c}
\varepsilon_{t} \\
0 \\
0
\end{array}\right]
$$

The eigenvalues are 1, 1, and 0.5, so $z_{t}$ has two unit roots (integrated of order 2: I(2) and needs to be differenced twice to become stationary).

Example 16 (Explosive $A R(1)$.) Consider the process $y_{t}=1.5 y_{t-1}+\varepsilon_{t}$. The eigenvalue is then outside the unit circle, so the process is explosive. This means that the impulse response to a shock to $\varepsilon_{t}$ diverges (it is $1.5^{s}$ for $s$ periods ahead).

### 2.7.2 Spurious Regressions

Strong trends often causes problems in econometric models where $y_{t}$ is regressed on $x_{t}$. In essence, if no trend is included in the regression, then $x_{t}$ will appear to be significant, just because it is a proxy for a trend. The same holds for unit root processes, even if they have no deterministic trends. However, the innovations accumulate and the series therefore tend to be trending in small samples. A warning sign of a spurious regression is when $R^{2}>D W$ statistics.

For trend-stationary data, this problem is easily solved by detrending with a linear trend (before estimating or just adding a trend to the regression).

However, this is usually a poor method for a unit root processes. What is needed is a first difference. For instance, a first difference of the random walk is

$$
\begin{align*}
\Delta y_{t} & =y_{t}-y_{t-1} \\
& =\varepsilon_{t}, \tag{2.49}
\end{align*}
$$

which is white noise (any finite difference, like $y_{t}-y_{t-s}$, will give a stationary series), so
we could proceed by applying standard econometric tools to $\Delta y_{t}$.
One may then be tempted to try first-differencing all non-stationary series, since it may be hard to tell if they are unit root process or just trend-stationary. For instance, a first difference of the trend stationary process, (2.46), gives

$$
\begin{equation*}
y_{t}-y_{t-1}=\beta+\varepsilon_{t}-\varepsilon_{t-1} . \tag{2.50}
\end{equation*}
$$

Its unclear if this is an improvement: the trend is gone, but the errors are now of MA(1) type (in fact, non-invertible, and therefore tricky, in particular for estimation).

### 2.7.3 Testing for a Unit Root I*

Suppose we run an OLS regression of

$$
\begin{equation*}
y_{t}=a y_{t-1}+\varepsilon_{t}, \tag{2.51}
\end{equation*}
$$

where the true value of $|a|<1$. The asymptotic distribution is of the LS estimator is

$$
\begin{equation*}
\sqrt{T}(\hat{a}-a) \sim N\left(0,1-a^{2}\right) . \tag{2.52}
\end{equation*}
$$

(The variance follows from the standard OLS formula where the variance of the estimator is $\sigma^{2}\left(X^{\prime} X / T\right)^{-1}$. Here plim $X^{\prime} X / T=\operatorname{Var}\left(y_{t}\right)$ which we know is $\left.\sigma^{2} /\left(1-a^{2}\right)\right)$.

It is well known (but not easy to show) that when $a=1$, then $\hat{a}$ is biased towards zero in small samples. In addition, the asymptotic distribution is no longer (2.52). In fact, there is a discontinuity in the limiting distribution as we move from a stationary/to a non-stationary variable. This, together with the small sample bias means that we have to use simulated critical values for testing the null hypothesis of $a=1$ based on the OLS estimate from (2.51).

The approach is to calculate the test statistic

$$
t=\frac{\hat{a}-1}{\operatorname{Std}(\hat{a})},
$$

and reject the null of non-stationarity if $t$ is less than the critical values published by Dickey and Fuller (typically more negative than the standard values to compensate for the small sample bias) or from your own simulations.

In principle, distinguishing between a stationary and a non-stationary series is very
difficult (and impossible unless we restrict the class of processes, for instance, to an $\operatorname{AR}(2)$ ), since any sample of a non-stationary process can be arbitrary well approximated by some stationary process et vice versa. The lesson to be learned, from a practical point of view, is that strong persistence in the data generating process (stationary or not) invalidates the usual results on inference. We are usually on safer ground to apply the unit root results in this case, even if the process is actually stationary.

### 2.7.4 Testing for a Unit Root II*

Reference: Fuller (1976), Introduction to Statistical Time Series; Dickey and Fuller (1979), "Distribution of the Estimators for Autoregressive Time Series with a Unit Root," Journal of the American Statistical Association, 74, 427-431.

Consider the $\operatorname{AR}(1)$ with intercept

$$
\begin{equation*}
y_{t}=\gamma+\alpha y_{t-1}+u_{t}, \text { or } \Delta y_{t}=\gamma+\beta y_{t-1}+u_{t}, \text { where } \beta=(\alpha-1) \tag{2.53}
\end{equation*}
$$

The $D F$ test is to test the null hypothesis that $\beta=0$, against $\beta<0$ using the usual $t$ statistic. However, under the null hypothesis, the distribution of the $t$ statistics is far from a student-t or normal distribution. Critical values, found in Fuller and Dickey and Fuller, are lower than the usual ones. Remember to add any nonstochastic regressors that in required, for instance, seasonal dummies, trends, etc. If you forget a trend, then the power of the test goes to zero as $T \rightarrow \infty$. The critical values are lower the more deterministic components that are added.

The asymptotic critical values are valid even under heteroskedasticity, and non-normal distributions of $u_{t}$. However, no autocorrelation in $u_{t}$ is allowed for. In contrast, the simulated small sample critical values are usually only valid for iid normally distributed disturbances.

The ADF test is a way to account for serial correlation in $u_{t}$. The same critical values apply. Consider an $\operatorname{AR}(1) u_{t}=\rho u_{t-1}+e_{t}$. A Cochrane-Orcutt transformation of (2.53) gives

$$
\begin{equation*}
\Delta y_{t}=\gamma(1-\rho)+\tilde{\beta} y_{t-1}+\rho(\beta+1) \Delta y_{t-1}+e_{t}, \text { where } \tilde{\beta}=\beta(1-\rho) . \tag{2.54}
\end{equation*}
$$

The test is here the $t$ test for $\tilde{\beta}$. The fact that $\tilde{\beta}=\beta(1-\rho)$ is of no importance, since $\tilde{\beta}$ is zero only if $\beta$ is (as long as $\rho<1$, as it must be). (2.54) generalizes so one should include
$p$ lags of $\Delta y_{t}$ if $u_{t}$ is an $\operatorname{AR}(p)$. The test remains valid even under an MA structure if the number of lags included increases at the rate $T^{1 / 3}$ as the sample lenngth increases. In practice: add lags until the remaining residual is white noise. The size of the test (probability of rejecting $\mathrm{H}_{0}$ when it is actually correct) can be awful in small samples for a series that is a $\mathrm{I}(1)$ process that initially "overshoots" over time, as $\Delta y_{t}=e_{t}-0.8 e_{t-1}$, since this makes the series look mean reverting (stationary). Similarly, the power (prob of rejecting $\mathrm{H}_{0}$ when it is false) can be awful when there is a lot of persistence, for instance, if $\alpha=0.95$.

The power of the test depends on the span of the data, rather than the number of observations. Seasonally adjusted data tend to look more integrated than they are. Should apply different critical values, see Ghysel and Perron (1993), Journal of Econometrics, 55, 57-98. A break in mean or trend also makes the data look non-stationary. Should perhaps apply tests that account for this, see Banerjee, Lumsdaine, Stock (1992), Journal of Business and Economics Statistics, 10, 271-287.

Park (1990, "Testing for Unit Roots and Cointegration by Variable Addition," Advances in Econometrics, 8, 107-133) sets up a framework where we can use both nonstationarity as the null hypothesis and where we can have stationarity as the null. Consider the regression

$$
\begin{equation*}
y_{t}=\sum_{s=0}^{p} \beta_{s} t^{s}+\sum_{s=p+1}^{q} \beta_{s} t^{s}+u_{t}, \tag{2.55}
\end{equation*}
$$

where the we want to test if $\mathrm{H}_{0}: \beta_{s}=0, s=p+1, \ldots, q$. If $F(p, q)$ is the Wald-statistics for this, then $J(p, q)=F(p, q) / T$ has some (complicated) asymptotic distribution under the null. You reject non-stationarity if $J(p, q)<$ critical value, since $J(p, q) \rightarrow^{p}$ 0 under (trend) stationarity.

Now, define

$$
\begin{equation*}
G(p, q)=F(p, q) \frac{\operatorname{Var}\left(u_{t}\right)}{\operatorname{Var}\left(\sqrt{T} \bar{u}_{t}\right)} \sim \chi_{p-q}^{2} \text { under } \mathrm{H}_{0} \text { of stationarity, } \tag{2.56}
\end{equation*}
$$

and $G(p, q) \rightarrow^{p} \infty$ under non-stationarity, so we reject stationarity if $G(p, q)>$ critical value. Note that $\operatorname{Var}\left(u_{t}\right)$ is a traditional variance, while $\operatorname{Var}\left(\sqrt{ } \bar{T} \bar{u}_{t}\right)$ can be estimated with a Newey-West estimator.

### 2.7.5 Cointegration*

Suppose $y_{1 t}$ and $y_{2 t}$ are both (scalar) unit root processes, but that

$$
\begin{align*}
z_{t} & =y_{1 t}-\beta y_{2 t}  \tag{2.57}\\
& =\left[\begin{array}{ll}
1 & -\beta
\end{array}\right]\left[\begin{array}{l}
y_{1 t} \\
y_{2 t}
\end{array}\right]
\end{align*}
$$

is stationary. The processes $y_{t}$ and $x_{t}$ must then share the same common stochastic trend, and are therefore cointegrated with the cointegrating vector $\left[\begin{array}{cc}1 & -\beta\end{array}\right]$. Running the regression (2.57) gives an estimator $\hat{\beta}_{L S}$ which converges much faster than usual (it is "superconsistent") and is not affected by any simultaneous equations bias. The intuition for the second result is that the simultaneous equations bias depends on the simultaneous reactions to the shocks, which are stationary and therefore without any long-run importance.

This can be generalized by letting $y_{t}$ be a vector of $n$ unit root processes which follows a VAR. For simplicity assume it is a $\operatorname{VAR}(2)$

$$
\begin{equation*}
y_{t}=A_{1} y_{t-1}+A_{2} y_{t-2}+\varepsilon_{t} . \tag{2.58}
\end{equation*}
$$

Subtract $y_{t}$ from both sides, add and subtract $A_{2} y_{t-1}$ from the right hand side

$$
\begin{align*}
y_{t}-y_{t-1} & =A_{1} y_{t-1}+A_{2} y_{t-2}+\varepsilon_{t}-y_{t-1}+A_{2} y_{t-1}-A_{2} y_{t-1} \\
& =\left(A_{1}+A_{2}-I\right) y_{t-1}-A_{2}\left(y_{t-1}-y_{t-2}\right)+\varepsilon_{t} \tag{2.59}
\end{align*}
$$

The left hand side is now stationary, and so is $y_{t-1}-y_{t-2}$ and $\varepsilon_{t}$ on the right hand side. It must therefore be the case that $\left(A_{1}+A_{2}-I\right) y_{t-1}$ is also stationary; it must be $n$ linear combinations of the cointegrating vectors. Since the number of cointegrating vectors must be less than $n$, the rank of $A_{1}+A_{2}-I$ must be less than $n$. To impose this calls for special estimation methods.

The simplest of these is Engle and Granger's two-step procedure. In the first step, we estimate the cointegrating vectors (as in 2.57) and calculate the different $z_{t}$ series (fewer than $n$ ). In the second step, these are used in the error correction form of the VAR

$$
\begin{equation*}
y_{t}-y_{t-1}=\gamma z_{t-1}-A_{2}\left(y_{t-1}-y_{t-2}\right)+\varepsilon_{t} \tag{2.60}
\end{equation*}
$$

to estimate $\gamma$ and $A_{2}$. The relation to (2.59) is most easily seen in the bivariate case. Then, by using (2.57) in (2.60) we get

$$
y_{t}-y_{t-1}=\left[\begin{array}{ll}
\gamma & -\gamma \beta \tag{2.61}
\end{array}\right] y_{t-1}-A_{2}\left(y_{t-1}-y_{t-2}\right)+\varepsilon_{t},
$$

so knowledge (estimates) of $\beta$ (scalar), $\gamma(2 \times 1), A_{2}(2 \times 2)$ allows us to "back out" $A_{1}$.

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## 3 The Distribution of a Sample Average

Reference: Hayashi (2000) 6.5
Additional references: Hamilton (1994) 14; Verbeek (2000) 4.10; Harris and Matyas (1999); and Pindyck and Rubinfeld (1997) Appendix 10.1; Cochrane (2001) 11.7

### 3.1 Variance of a Sample Average

In order to understand the distribution of many estimators we need to get an important building block: the variance of a sample average.

Consider a covariance stationary vector process $m_{t}$ with zero mean and $\operatorname{Cov}\left(m_{t}, m_{t-s}\right)=$ $R(s)$ (which only depends on $s$ ). That is, we allow for serial correlation in $m_{t}$, but no heteroskedasticity. This is more restrictive than we want, but we will return to that further on.

Let $\bar{m}=\sum_{t=1}^{T} m_{t} / T$. The sampling variance of a mean estimator of the zero mean random variable $m_{t}$ is defined as

$$
\begin{equation*}
\operatorname{Cov}(\bar{m})=\mathrm{E}\left[\left(\frac{1}{T} \sum_{t=1}^{T} m_{t}\right)\left(\frac{1}{T} \sum_{\tau=1}^{T} m_{\tau}\right)^{\prime}\right] \tag{3.1}
\end{equation*}
$$

Let the covariance (matrix) at lag $s$ be

$$
\begin{align*}
R(s) & =\operatorname{Cov}\left(m_{t}, m_{t-s}\right) \\
& =\mathrm{E} m_{t} m_{t-s}^{\prime}, \tag{3.2}
\end{align*}
$$

since $\mathrm{E} m_{t}=0$ for all $t$.

Example 1 ( $m_{t}$ is a scalar iid process.) When $m_{t}$ is a scalar iid process, then

$$
\begin{aligned}
\operatorname{Var}\left(\frac{1}{T} \sum_{t=1}^{T} m_{t}\right) & =\frac{1}{T^{2}} \sum_{t=1}^{T} \operatorname{Var}\left(m_{t}\right) / * \text { independently distributed } * / \\
& =\frac{1}{T^{2}} T \operatorname{Var}\left(m_{t}\right) / * \text { identically distributed } * / \\
& =\frac{1}{T} \operatorname{Var}\left(m_{t}\right)
\end{aligned}
$$

This is the classical iid case. Clearly, $\lim _{T \Rightarrow \infty} \operatorname{Var}(\bar{m})=0$. By multiplying both sides by $T$ we instead get $\operatorname{Var}(\sqrt{T} \bar{m})=\operatorname{Var}\left(m_{t}\right)$, which is often more convenient for asymptotics.

Example 2 Let $x_{t}$ and $z_{t}$ be two scalars, with samples averages $\bar{x}$ and $\bar{z}$. Let $m_{t}=$ $\left[\begin{array}{ll}x_{t} & z_{t}\end{array}\right]^{\prime}$. Then $\operatorname{Cov}(\bar{m})$ is

$$
\begin{aligned}
\operatorname{Cov}(\bar{m}) & =\operatorname{Cov}\left(\left[\begin{array}{l}
\bar{x} \\
\bar{z}
\end{array}\right]\right) \\
& =\left[\begin{array}{cc}
\operatorname{Var}(\bar{x}) & \operatorname{Cov}(\bar{x}, \bar{z}) \\
\operatorname{Cov}(\bar{z}, \bar{x}) & \operatorname{Var}(\bar{z})
\end{array}\right] .
\end{aligned}
$$

Example $3(\operatorname{Cov}(\bar{m})$ with $T=3$.) With $T=3$, we have

$$
\operatorname{Cov}(T \bar{m})=
$$

$E\left(m_{1}+m_{2}+m_{3}\right)\left(m_{1}^{\prime}+m_{2}^{\prime}+m_{3}^{\prime}\right)=$
$\underbrace{E\left(m_{1} m_{1}^{\prime}+m_{2} m_{2}^{\prime}+m_{3} m_{3}^{\prime}\right)}_{3 R(0)}+\underbrace{E\left(m_{2} m_{1}^{\prime}+m_{3} m_{2}^{\prime}\right)}_{2 R(1)}+\underbrace{E\left(m_{1} m_{2}^{\prime}+m_{2} m_{3}^{\prime}\right)}_{2 R(-1)}+\underbrace{E m_{3} m_{1}^{\prime}}_{R(2)}+\underbrace{E m_{1} m_{3}^{\prime}}_{R(-2)}$.
The general pattern in the previous example is

$$
\begin{equation*}
\operatorname{Cov}(T \bar{m})=\sum_{s=-(T-1)}^{T-1}(T-|s|) R(s) \tag{3.3}
\end{equation*}
$$

Divide both sides by $T$

$$
\begin{equation*}
\operatorname{Cov}(\sqrt{T} \bar{m})=\sum_{s=-(T-1)}^{T-1}\left(1-\frac{|s|}{T}\right) R(s) . \tag{3.4}
\end{equation*}
$$

This is the exact expression for a given sample size.

In many cases, we use the asymptotic expression (limiting value as $T \rightarrow \infty$ ) instead. If $R(s)=0$ for $s>q$ so $m_{t}$ is an $\operatorname{MA}(q)$, then the limit as the sample size goes to infinity is

$$
\begin{equation*}
\operatorname{ACov}(\sqrt{T} \bar{m})=\lim _{T \rightarrow \infty} \operatorname{Cov}(\sqrt{T} \bar{m})=\sum_{s=-q}^{q} R(s) \tag{3.5}
\end{equation*}
$$

where ACov stands for the asymptotic variance-covariance matrix. This continues to hold even if $q=\infty$, provided $R(s)$ goes to zero sufficiently quickly, as it does in stationary VAR systems. In this case we have

$$
\begin{equation*}
\operatorname{ACov}(\sqrt{T} \bar{m})=\sum_{s=-\infty}^{\infty} R(s) \tag{3.6}
\end{equation*}
$$

Estimation in finite samples will of course require some cut-off point, which is discussed below.

The traditional estimator of $\operatorname{ACov}(\sqrt{T} \bar{m})$ is just $R(0)$, which is correct when $m_{t}$ has no autocorrelation, that is

$$
\begin{equation*}
\operatorname{ACov}(\sqrt{T} \bar{m})=R(0)=\operatorname{Cov}\left(m_{t}, m_{t}\right) \text { if } \operatorname{Cov}\left(m_{t}, m_{t-s}\right) \text { for } s \neq 0 \tag{3.7}
\end{equation*}
$$

By comparing with (3.5) we see that this underestimates the true variance of autocovariances are mostly positive, and overestimates if they are mostly negative. The errors can be substantial.

Example 4 (Variance of sample mean of $A R(1)$.) Let $m_{t}=\rho m_{t-1}+u_{t}$, where $\operatorname{Var}\left(u_{t}\right)=$ $\sigma^{2}$. Note that $R(s)=\rho^{|s|} \sigma^{2} /\left(1-\rho^{2}\right)$, so

$$
\begin{aligned}
\operatorname{AVar}(\sqrt{T} \bar{m}) & =\sum_{s=-\infty}^{\infty} R(s) \\
& =\frac{\sigma^{2}}{1-\rho^{2}} \sum_{s=-\infty}^{\infty} \rho^{|s|}=\frac{\sigma^{2}}{1-\rho^{2}}\left(1+2 \sum_{s=1}^{\infty} \rho^{s}\right) \\
& =\frac{\sigma^{2}}{1-\rho^{2}} \frac{1+\rho}{1-\rho}
\end{aligned}
$$

which is increasing in $\rho$ (provided $|\rho|<1$, as required for stationarity). The variance of $\bar{m}$ is much larger for $\rho$ close to one than for $\rho$ close to zero: the high autocorrelation create long swings, so the mean cannot be estimated with any good precision in a small


Figure 3.1: Variance of $\sqrt{T}$ times sample mean of $\operatorname{AR}(1)$ process $m_{t}=\rho m_{t-1}+u_{t}$.
sample. If we disregard all autocovariances, then we would conclude that the variance of $\sqrt{T} \bar{m}$ is $\sigma^{2} /\left(1-\rho^{2}\right)$, which is smaller (larger) than the true value when $\rho>0(\rho<0)$. For instance, with $\rho=0.85$, it is approximately 12 times too small. See Figure 3.1.a for an illustration.

Example 5 (Variance of sample mean of $A R(1)$, continued.) Part of the reason why $\operatorname{Var}(\bar{m})$ increased with $\rho$ in the previous examples is that $\operatorname{Var}\left(m_{t}\right)$ increases with $\rho$. We can eliminate this effect by considering how much larger $A \operatorname{Var}(\sqrt{T} \bar{m})$ is than in the iid case, that is, $A \operatorname{Var}(\sqrt{T} \bar{m}) / \operatorname{Var}\left(m_{t}\right)=(1+\rho) /(1-\rho)$. This ratio is one for $\rho=0$ (iid data), less than one for $\rho<0$, and greater than one for $\pi>0$. This says that if relatively more of the variance in $m_{t}$ comes from long swings (high $\rho$ ), then the sample mean is more uncertain. See Figure 3.1.b for an illustration.

Example 6 (Variance of sample mean of $A R(1)$, illustration of why $\lim _{T \rightarrow \infty}$ of (3.4).) For an $\operatorname{AR}(1)$ (3.4) is

$$
\begin{aligned}
\operatorname{Var}(\sqrt{T} \bar{m}) & =\frac{\sigma^{2}}{1-\rho^{2}} \sum_{s=-(T-1)}^{T-1}\left(1-\frac{|s|}{T}\right) \rho^{|s|} \\
& =\frac{\sigma^{2}}{1-\rho^{2}}\left[1+2 \sum_{s=1}^{T-1}\left(1-\frac{s}{T}\right) \rho^{s}\right] \\
& =\frac{\sigma^{2}}{1-\rho^{2}}\left[1+2 \frac{\rho}{1-\rho}+2 \frac{\rho^{T+1}-\rho}{T(1-\rho)^{2}}\right] .
\end{aligned}
$$

The last term in brackets goes to zero as $T$ goes to infinity. We then get the result in Example 4.

### 3.2 The Newey-West Estimator

### 3.2.1 Definition of the Estimator

Newey and West (1987) suggested the following estimator of the covariance matrix in (3.5) as (for some $n<T$ )

$$
\begin{align*}
\widehat{\operatorname{ACov}}(\sqrt{T} \bar{m})= & \sum_{s=-n}^{n}\left(1-\frac{|s|}{n+1}\right) \hat{R}(s) \\
= & \hat{R}(0)+\sum_{s=1}^{n}\left(1-\frac{s}{n+1}\right)(\hat{R}(s)+\hat{R}(-s)), \text { or since } \hat{R}(-s)=\hat{R}^{\prime}(s) \\
= & \hat{R}(0)+\sum_{s=1}^{n}\left(1-\frac{s}{n+1}\right)\left(\hat{R}(s)+\hat{R}^{\prime}(s)\right), \text { where }  \tag{3.8}\\
& \hat{R}(s)=\frac{1}{T} \sum_{t=s+1}^{T} m_{t} m_{t-s}^{\prime}\left(\text { if } \mathrm{E} m_{t}=0\right) . \tag{3.9}
\end{align*}
$$

The tent shaped (Bartlett) weights in (3.8) guarantee a positive definite covariance estimate. In contrast, equal weights (as in (3.5)), may give an estimated covariance matrix which is not positive definite, which is fairly awkward. Newey and West (1987) showed that this estimator is consistent if we let $n$ go to infinity as $T$ does, but in such a way that $n / T^{1 / 4}$ goes to zero.

There are several other possible estimators of the covariance matrix in (3.5), but simulation evidence suggest that they typically do not improve a lot on the Newey-West estimator.

Example 7 ( $m_{t}$ is $M A(1)$ ).) Suppose we know that $m_{t}=\varepsilon_{t}+\theta \varepsilon_{t-1}$. Then $R(s)=0$ for $s \geq 2$, so it might be tempting to use $n=1$ in (3.8). This gives $\widehat{A C o v}(\sqrt{T} \bar{m})=\hat{R}(0)+$ $\frac{1}{2}\left[\hat{R}(1)+\hat{R}^{\prime}(1)\right]$, while the theoretical expression (3.5) is ACov $=R(0)+R(1)+R^{\prime}(1)$. The Newey-West estimator puts too low weights on the first lead and lag, which suggests that we should use $n>1$ (or more generally, $n>q$ for an $M A(q)$ process).

It can also be shown that, under quite general circumstances, $\hat{S}$ in (3.8)-(3.9) is a consistent estimator of $A \operatorname{Cov}(\sqrt{T} \bar{m})$, even if $m_{t}$ is heteroskedastic (on top of being autocorrelated). (See Hamilton (1994) 10.5 for a discussion.)

### 3.2.2 How to Implement the Newey-West Estimator

Economic theory and/or stylized facts can sometimes help us choose the lag length $n$. For instance, we may have a model of stock returns which typically show little autocorrelation, so it may make sense to set $n=0$ or $n=1$ in that case. A popular choice of $n$ is to round $(T / 100)^{1 / 4}$ down to the closest integer, although this does not satisfy the consistency requirement.

It is important to note that definition of the covariance matrices in (3.2) and (3.9) assume that $m_{t}$ has zero mean. If that is not the case, then the mean should be removed in the calculation of the covariance matrix. In practice, you remove the same number, estimated on the whole sample, from both $m_{t}$ and $m_{t-s}$. It is often recommended to remove the sample means even if theory tells you that the true mean is zero.

### 3.3 Summary

$$
\begin{aligned}
\text { Let } \bar{m} & =\frac{1}{T} \sum_{t=1}^{T} m_{t} \text { and } R(s)=\operatorname{Cov}\left(m_{t}, m_{t-s}\right) . \text { Then } \\
\text { ACov }(\sqrt{T} \bar{m}) & =\sum_{s=-\infty}^{\infty} R(s) \\
\text { ACov }(\sqrt{T} \bar{m}) & =R(0)=\operatorname{Cov}\left(m_{t}, m_{t}\right) \text { if } R(s)=\mathbf{0} \text { for } s \neq 0 \\
\text { Newey-West }: & \widehat{\operatorname{ACov}}(\sqrt{T} \bar{m})=\hat{R}(0)+\sum_{s=1}^{n}\left(1-\frac{s}{n+1}\right)\left(\hat{R}(s)+\hat{R}^{\prime}(s)\right) .
\end{aligned}
$$

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## 4 Least Squares

Reference: Greene (2000) 6
Additional references: Hayashi (2000) 1-2; Verbeek (2000) 1-4; Hamilton (1994) 8

### 4.1 Definition of the LS Estimator

### 4.1.1 LS with Summation Operators

Consider the linear model

$$
\begin{equation*}
y_{t}=x_{t}^{\prime} \beta_{0}+u_{t} \tag{4.1}
\end{equation*}
$$

where $y_{t}$ and $u_{t}$ are scalars, $x_{t}$ a $k \times 1$ vector, and $\beta_{0}$ is a $k \times 1$ vector of the true coefficients. Least squares minimizes the sum of the squared fitted residuals

$$
\begin{equation*}
\sum_{t=1}^{T} e_{t}^{2}=\sum_{t=1}^{T}\left(y_{t}-x_{t}^{\prime} \beta\right)^{2} \tag{4.2}
\end{equation*}
$$

by choosing the vector $\beta$. The first order conditions are

$$
\begin{align*}
\mathbf{0}_{k x 1} & =\sum_{t=1}^{T} x_{t}\left(y_{t}-x_{t}^{\prime} \hat{\beta}_{L S}\right) \text { or }  \tag{4.3}\\
\sum_{t=1}^{T} x_{t} y_{t} & =\sum_{t=1}^{T} x_{t} x_{t}^{\prime} \hat{\beta}_{L S} \tag{4.4}
\end{align*}
$$

which are the so called normal equations. These can be solved as

$$
\begin{align*}
\hat{\beta}_{L S} & =\left(\sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right)^{-1} \sum_{t=1}^{T} x_{t} y_{t}  \tag{4.5}\\
& =\left(\frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} x_{t} y_{t} \tag{4.6}
\end{align*}
$$

Remark 1 (Summation and vectors) Let $z_{t}$ and $x_{t}$ be the vectors

$$
z_{t}=\left[\begin{array}{l}
z_{1 t} \\
z_{2 t}
\end{array}\right] \text { and } x_{t}=\left[\begin{array}{l}
x_{1 t} \\
x_{2 t} \\
x_{3 t}
\end{array}\right]
$$

then
$\sum_{t=1}^{T} x_{t} z_{t}^{\prime}=\sum_{t=1}^{T}\left[\begin{array}{c}x_{1 t} \\ x_{2 t} \\ x_{3 t}\end{array}\right]\left[\begin{array}{ll}z_{1 t} & z_{2 t}\end{array}\right]=\sum_{t=1}^{T}\left[\begin{array}{cc}x_{1 t} z_{1 t} & x_{1 t} z_{2 t} \\ x_{2 t} z_{1 t} & x_{2 t} z_{2 t} \\ x_{3 t} z_{1 t} & x_{3 t} z_{2 t}\end{array}\right]=\left[\begin{array}{cc}\sum_{t=1}^{T} x_{1 t} z_{1 t} & \sum_{t=1}^{T} x_{1 t} z_{2 t} \\ \sum_{t=1}^{T} x_{2 t} z_{1 t} & \sum_{t=1}^{T} x_{2 t} z_{2 t} \\ \sum_{t=1}^{T} x_{3 t} z_{1 t} & \sum_{t=1}^{T} x_{3 t} z_{2 t}\end{array}\right]$

### 4.1.2 LS in Matrix Form

Define the matrices

$$
Y=\left[\begin{array}{c}
y_{1}  \tag{4.7}\\
y_{2} \\
\vdots \\
y_{T}
\end{array}\right]_{T \times 1}, u=\left[\begin{array}{c}
u_{1} \\
u_{2} \\
\vdots \\
u_{T}
\end{array}\right]_{T \times 1}, X=\left[\begin{array}{c}
x_{1}^{\prime} \\
x_{2}^{\prime} \\
\vdots \\
x_{T}^{\prime}
\end{array}\right]_{T \times k}, \text { and } e=\left[\begin{array}{c}
e_{1} \\
e_{2} \\
\vdots \\
e_{T}
\end{array}\right]_{T \times 1} .
$$

Write the model (4.1) as

$$
\begin{align*}
{\left[\begin{array}{c}
y_{1} \\
y_{2} \\
\vdots \\
y_{T}
\end{array}\right] } & =\left[\begin{array}{c}
x_{1}^{\prime} \\
x_{2}^{\prime} \\
\vdots \\
x_{T}^{\prime}
\end{array}\right] \beta_{0}+\left[\begin{array}{c}
u_{1} \\
u_{2} \\
\vdots \\
u_{T}
\end{array}\right] \text { or }  \tag{4.8}\\
Y & =X \beta_{0}+u . \tag{4.9}
\end{align*}
$$

Remark 2 Let $x_{t}$ be a $k \times 1$ and $z_{t}$ an $m \times 1$ vector. Define the matrices

$$
X=\left[\begin{array}{c}
x_{1}^{\prime} \\
x_{2}^{\prime} \\
\vdots \\
x_{T}^{\prime}
\end{array}\right]_{T \times k} \text { and } Z=\left[\begin{array}{c}
z_{1}^{\prime} \\
z_{2}^{\prime} \\
\vdots \\
z_{T}^{\prime}
\end{array}\right]_{T \times m}
$$

We then have

$$
\sum_{t=1}^{T} x_{t} z_{t}^{\prime}=X^{\prime} Z
$$

We can then rewrite the loss function (4.2) as $e^{\prime} e$, the first order conditions (4.3) and (4.4) as (recall that $y_{t}=y_{t}^{\prime}$ since it is a scalar)

$$
\begin{align*}
& \mathbf{0}_{k x 1}=X^{\prime}\left(Y-X \hat{\beta}_{L S}\right)  \tag{4.10}\\
& X^{\prime} Y=X^{\prime} X \hat{\beta}_{L S}, \tag{4.11}
\end{align*}
$$

and the solution (4.5) as

$$
\begin{equation*}
\hat{\beta}_{L S}=\left(X^{\prime} X\right)^{-1} X^{\prime} Y \tag{4.12}
\end{equation*}
$$

### 4.2 LS and $R^{2}$ *

The first order conditions in LS are

$$
\begin{equation*}
\sum_{t=1}^{T} x_{t} \hat{u}_{t}=\mathbf{0}, \text { where } \hat{u}_{t}=y_{t}-\hat{y}_{t}, \text { with } \hat{y}_{t}=x_{t}^{\prime} \hat{\beta} \tag{4.13}
\end{equation*}
$$

This implies that the fitted residuals and fitted values are orthogonal, $\Sigma_{t=1}^{T} \hat{y}_{t} \hat{u}_{t}=\Sigma_{t=1}^{T} \hat{\beta}^{\prime} x_{t} \hat{u}_{t}=$ 0 . If we let $x_{t}$ include a constant, then (4.13) also implies that the fitted residuals have a zero mean, $\Sigma_{t=1}^{T} \hat{u}_{t} / T=0$. We can then decompose the sample variance (denoted $\widehat{\operatorname{Var}}$ ) of $y_{t}=\hat{y}_{t}+\hat{u}_{t}$ as

$$
\begin{equation*}
\widehat{\operatorname{Var}}\left(y_{t}\right)=\widehat{\operatorname{Var}}\left(\hat{y}_{t}\right)+\widehat{\operatorname{Var}}\left(\hat{u}_{t}\right) \tag{4.14}
\end{equation*}
$$

since $\hat{y}_{t}$ and $\hat{u}_{t}$ are uncorrelated in this case. (Note that $\operatorname{Cov}\left(\hat{y}_{t}, \hat{u}_{t}\right)=\mathrm{E} \hat{y}_{t} \hat{u}_{t}-\mathrm{E} \hat{y}_{t} \mathrm{E} \hat{u}_{t}$ so the orthogonality is not enough to allow the decomposition; we also need $\mathrm{E} \hat{y}_{t} \mathrm{E} \hat{u}_{t}=0$ this holds for sample moments as well.)

We define $R^{2}$ as the fraction of $\widehat{\operatorname{Var}}\left(y_{t}\right)$ that is explained by the model

$$
\begin{align*}
R^{2} & =\frac{\widehat{\operatorname{Var}}\left(\hat{y}_{t}\right)}{\widehat{\operatorname{Var}}\left(y_{t}\right)}  \tag{4.15}\\
& =1-\frac{\widehat{\operatorname{Var}}\left(\hat{u}_{t}\right)}{\widehat{\operatorname{Var}}\left(y_{t}\right)} . \tag{4.16}
\end{align*}
$$

LS minimizes the sum of squared fitted errors, which is proportional to $\widehat{\operatorname{Var}}\left(\hat{u}_{t}\right)$, so it maximizes $R^{2}$.

We can rewrite $R^{2}$ by noting that

$$
\begin{equation*}
\widehat{\operatorname{Cov}}\left(y_{t}, \hat{y}_{t}\right)=\widehat{\operatorname{Cov}}\left(\hat{y}_{t}+\hat{u}_{t}, \hat{y}_{t}\right)=\widehat{\operatorname{Tar}}\left(\hat{y}_{t}\right) . \tag{4.17}
\end{equation*}
$$

Use this to substitute for $\widehat{\operatorname{Var}}\left(\hat{y}_{t}\right)$ in (4.15) and then multiply both sides with $\widehat{\operatorname{Cov}}\left(y_{t}, \hat{y}_{t}\right) / \widehat{\operatorname{Var}}\left(\hat{y}_{t}\right)=$ 1 to get

$$
\begin{align*}
R^{2} & =\frac{\widehat{\operatorname{Cov}}\left(y_{t}, \hat{y}_{t}\right)^{2}}{\widehat{\operatorname{Var}}\left(y_{t}\right) \widehat{\operatorname{Tar}}\left(\hat{y}_{t}\right)} \\
& =\widehat{\operatorname{Corr}}\left(y_{t}, \hat{y}_{t}\right)^{2} \tag{4.18}
\end{align*}
$$

which shows that $R^{2}$ is the square of correlation coefficient of the actual and fitted value. Note that this interpretation of $R^{2}$ relies on the fact that $\widehat{\operatorname{Cov}}\left(\hat{y}_{t}, \hat{u}_{t}\right)=0$. From (4.14) this implies that the sample variance of the fitted variables is smaller than the sample variance of $y_{t}$. From (4.15) we see that this implies that $0 \leq R^{2} \leq 1$.

To get a bit more intuition for what $R^{2}$ represents, suppose the estimated coefficients equal the true coefficients, so $\hat{y}_{t}=x_{t}^{\prime} \beta_{0}$. In this case, $R^{2}=\operatorname{Corr}\left(x_{t}^{\prime} \beta_{0}+u_{t}, x_{t}^{\prime} \beta_{0}\right)^{2}$, that is, the squared correlation of $y_{t}$ with the systematic part of $y_{t}$. Clearly, if the model is perfect so $u_{t}=0$, then $R^{2}=1$. On contrast, when there is no movements in the systematic part ( $\beta_{0}=\mathbf{0}$ ), then $R^{2}=0$.

Remark 3 In a simple regression where $y_{t}=a+b x_{t}+u_{t}$, where $x_{t}$ is a scalar, $R^{2}=$ $\widehat{\operatorname{Corr}}\left(y_{t}, x_{t}\right)^{2}$. To see this, note that, in this case (4.18) can be written

$$
R^{2}=\frac{\widehat{\operatorname{Cov}}\left(y_{t}, \hat{b} x_{t}\right)^{2}}{\widehat{\operatorname{Var}}\left(y_{t}\right) \widehat{\operatorname{Var}}\left(\hat{b} x_{t}\right)}=\frac{\hat{b}^{2} \widehat{\operatorname{Cov}}\left(y_{t}, x_{t}\right)^{2}}{\hat{b}^{2} \widehat{\operatorname{Var}}\left(y_{t}\right) \widehat{\operatorname{Var}}\left(x_{t}\right)},
$$

so the $\hat{b}^{2}$ terms cancel.
Remark 4 Now, consider the reverse regression $x_{t}=c+d y_{t}+v_{t}$. The LS estimator of the slope is $\hat{d}_{L S}=\widehat{\operatorname{Cov}}\left(y_{t}, x_{t}\right) / \widehat{\operatorname{Var}}\left(y_{t}\right)$. Recall that $\hat{b}_{L S}=\widehat{\operatorname{Cov}}\left(y_{t}, x_{t}\right) / \widehat{\operatorname{Var}}\left(x_{t}\right)$. We therefore have

$$
\hat{b}_{L S} \hat{d}_{L S}=\frac{\widehat{\operatorname{Cov}}\left(y_{t}, x_{t}\right)^{2}}{\widehat{\operatorname{Var}}\left(y_{t}\right) \widehat{\operatorname{Var}}\left(x_{t}\right)}=R^{2}
$$

This shows that $\hat{d}_{L S}=1 / \hat{b}_{L S}$ if (and only if) $R^{2}=1$.

### 4.3 Finite Sample Properties of LS

Use the true model (4.1) to substitute for $y_{t}$ in the definition of the LS estimator (4.6)

$$
\begin{align*}
\hat{\beta}_{L S} & =\left(\frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} x_{t}\left(x_{t}^{\prime} \beta_{0}+u_{t}\right) \\
& =\beta_{0}+\left(\frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} x_{t} u_{t} \tag{4.19}
\end{align*}
$$

It is possible to show unbiasedness of the LS estimator, even if $x_{t}$ stochastic and $u_{t}$ is autocorrelated and heteroskedastic-provided $\mathrm{E}\left(u_{t} \mid x_{t-s}\right)=0$ for all $s$. Let $\mathrm{E}\left(u_{t} \mid\left\{x_{t}\right\}_{t=1}^{T}\right)$ denote the expectation of $u_{t}$ conditional on all values of $x_{t-s}$. Using iterated expectations on (4.19) then gives

$$
\begin{align*}
\mathrm{E} \hat{\beta}_{L S} & =\beta_{0}+\mathrm{E}_{x}\left[\left(\frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} x_{t} \mathrm{E}\left(u_{t} \mid\left\{x_{t}\right\}_{t=1}^{T}\right)\right]  \tag{4.20}\\
& =\beta_{0} \tag{4.21}
\end{align*}
$$

since $\mathrm{E}\left(u_{t} \mid x_{t-s}\right)=0$ for all $s$. This is, for instance, the case when the regressors are deterministic. Notice that $\mathrm{E}\left(u_{t} \mid x_{t}\right)=0$ is not enough for unbiasedness since (4.19) contains terms involving $x_{t-s} x_{t} u_{t}$ from the product of $\left(\frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right)^{-1}$ and $x_{t} u_{t}$.

Example 5 (AR(1).) Consider estimating $\alpha$ in $y_{t}=\alpha y_{t-1}+u_{t}$. The LS estimator is

$$
\begin{aligned}
\hat{\alpha}_{L S} & =\left(\frac{1}{T} \sum_{t=1}^{T} y_{t-1}^{2}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} y_{t-1} y_{t} \\
& =\alpha+\left(\frac{1}{T} \sum_{t=1}^{T} y_{t-1}^{2}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} y_{t-1} u_{t} .
\end{aligned}
$$

In this case, the assumption $E\left(u_{t} \mid x_{t-s}\right)=0$ for all $s$ (that is, $\left.s=\ldots,-1,0,1, \ldots\right)$ is false, since $x_{t+1}=y_{t}$ and $u_{t}$ and $y_{t}$ are correlated. We can therefore not use this way of proving that $\hat{\alpha}_{L S}$ is unbiased. In fact, it is not, and it can be shown that $\hat{\alpha}_{L S}$ is downward-biased

## if $\alpha>0$, and that this bias gets quite severe as $\alpha$ gets close to unity.

The finite sample distribution of the LS estimator is typically unknown.
Even in the most restrictive case where $u_{t}$ is iid $N\left(0, \sigma^{2}\right)$ and $\mathrm{E}\left(u_{t} \mid x_{t-s}\right)=0$ for all $s$, we can only get that

$$
\begin{equation*}
\hat{\beta}_{L S} \left\lvert\,\left\{x_{t}\right\}_{t=1}^{T} \sim N\left[\beta_{0}, \sigma^{2}\left(\frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right)^{-1}\right]\right. \tag{4.22}
\end{equation*}
$$

This says that the estimator, conditional on the sample of regressors, is normally distributed. With deterministic $x_{t}$, this clearly means that $\hat{\beta}_{L S}$ is normally distributed in a small sample. The intuition is that the LS estimator with deterministic regressors is just a linear combination of the normally distributed $y_{t}$, so it must be normally distributed. However, if $x_{t}$ is stochastic, then we have to take into account the distribution of $\left\{x_{t}\right\}_{t=1}^{T}$ to find the unconditional distribution of $\hat{\beta}_{L S}$. The principle is that

$$
\operatorname{pdf}(\hat{\beta})=\int_{-\infty}^{\infty} \operatorname{pdf}(\hat{\beta}, x) d x=\int_{-\infty}^{\infty} \operatorname{pdf}(\hat{\beta} \mid x) \operatorname{pdf}(x) d x
$$

so the distribution in (4.22) must be multiplied with the probability density function of $\left\{x_{t}\right\}_{t=1}^{T}$ and then integrated over $\left\{x_{t}\right\}_{t=1}^{T}$ to give the unconditional distribution (marginal) of $\hat{\beta}_{L S}$. This is typically not a normal distribution.

Another way to see the same problem is to note that $\hat{\beta}_{L S}$ in (4.19) is a product of two random variables, $\left(\Sigma_{t=1}^{T} x_{t} x_{t}^{\prime} / T\right)^{-1}$ and $\Sigma_{t=1}^{T} x_{t} u_{t} / T$. Even if $u_{t}$ happened to be normally distributed, there is no particular reason why $x_{t} u_{t}$ should be, and certainly no strong reason for why $\left(\Sigma_{t=1}^{T} x_{t} x_{t}^{\prime} / T\right)^{-1} \Sigma_{t=1}^{T} x_{t} u_{t} / T$ should be.

### 4.4 Consistency of LS

Reference: Greene (2000) 9.3-5 and 11.2; Hamilton (1994) 8.2; Davidson (2000) 3 We now study if the LS estimator is consistent.
Remark 6 Suppose the true parameter value is $\beta_{0}$. The estimator $\hat{\beta}_{T}$ (which, of course, depends on the sample size $T$ ) is said to be consistent if for every $\varepsilon>0$ and $\delta>0$ there exists $N$ such that for $T \geq N$

$$
\operatorname{Pr}\left(\left\|\hat{\beta}_{T}-\beta_{0}\right\|>\delta\right)<\varepsilon
$$

$\left(\|x\|=\sqrt{x^{\prime} x}\right.$, the Euclidean distance of $x$ from zero. $)$ We write this $\operatorname{plim} \hat{\beta}_{T}=\beta_{0}$ or just $\operatorname{plim} \hat{\beta}=\beta_{0}$, or perhaps $\hat{\beta} \rightarrow^{p} \beta_{0}$. (For an estimator of a covariance matrix, the most convenient is to stack the unique elements in a vector and then apply the definition above.)

Remark 7 (Slutsky's theorem.) If $g($.$) is a continuous function, then \operatorname{plim} g\left(z_{T}\right)=$ $g\left(\operatorname{plim} z_{T}\right)$. In contrast, note that $E g\left(z_{T}\right)$ is generally not equal to $g\left(E z_{T}\right)$, unless $g($. is a linear function.

Remark 8 (Probability limit of product.) Let $x_{T}$ and $y_{T}$ be two functions of a sample of length $T$. If $\operatorname{plim} x_{T}=a$ and $\operatorname{plim} y_{T}=b$, then $\operatorname{plim} x_{T} y_{T}=a b$.

Assume

$$
\begin{equation*}
\operatorname{plim} \frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}=\Sigma_{x x}<\infty, \text { and } \Sigma_{x x} \text { invertible. } \tag{4.23}
\end{equation*}
$$

The plim carries over to the inverse by Slutsky's theorem. ${ }^{1}$ Use the facts above to write the probability limit of (4.19) as

$$
\begin{equation*}
\operatorname{plim} \hat{\beta}_{L S}=\beta_{0}+\Sigma_{x x}^{-1} \operatorname{plim} \frac{1}{T} \sum_{t=1}^{T} x_{t} u_{t} \tag{4.24}
\end{equation*}
$$

To prove consistency of $\hat{\beta}_{L S}$ we therefore have to show that

$$
\begin{equation*}
\operatorname{plim} \frac{1}{T} \sum_{t=1}^{T} x_{t} u_{t}=\mathrm{E} x_{t} u_{t}=\operatorname{Cov}\left(x_{t}, u_{t}\right)=0 \tag{4.25}
\end{equation*}
$$

This is fairly easy to establish in special cases, for instance, when $w_{t}=x_{t} u_{t}$ is iid or when there is either heteroskedasticity or serial correlation. The case with both serial correlation and heteroskedasticity is just a bit more complicated. In other cases, it is clear that the covariance the residuals and the regressors are not all zero-for instance when some of the regressors are measured with error or when some of them are endogenous variables.

[^1]An example of a case where LS is not consistent is when the errors are autocorrelated and the regressors include lags of the dependent variable. For instance, suppose the error is a $\mathrm{MA}(1)$ process

$$
\begin{equation*}
u_{t}=\varepsilon_{t}+\theta_{1} \varepsilon_{t-1}, \tag{4.26}
\end{equation*}
$$

where $\varepsilon_{t}$ is white noise and that the regression equation is an $\operatorname{AR}(1)$

$$
\begin{equation*}
y_{t}=\rho y_{t-1}+u_{t} . \tag{4.27}
\end{equation*}
$$

This is an $\operatorname{ARMA}(1,1)$ model and it is clear that the regressor and error in (4.27) are correlated, so LS is not a consistent estimator of an $\operatorname{ARMA}(1,1)$ model.

### 4.5 Asymptotic Normality of LS

Reference: Greene (2000) 9.3-5 and 11.2; Hamilton (1994) 8.2; Davidson (2000) 3
Remark 9 (Continuous mapping theorem.) Let the sequences of random matrices $\left\{x_{T}\right\}$ and $\left\{y_{T}\right\}$, and the non-random matrix $\left\{a_{T}\right\}$ be such that $x_{T} \xrightarrow{d} x, y_{T} \xrightarrow{p} y$, and $a_{T} \rightarrow a$ (a traditional limit). Let $g\left(x_{T}, y_{T}, a_{T}\right)$ be a continuous function. Then $g\left(x_{T}, y_{T}, a_{T}\right) \xrightarrow{d}$ $g(x, y, a)$. Either of $y_{T}$ and $a_{T}$ could be irrelevant in $g$.

Remark 10 From the previous remark: if $x_{T} \xrightarrow{d} x$ (a random variable) and $\operatorname{plim} Q_{T}=$ $Q$ (a constant matrix), then $Q_{T} x_{T} \xrightarrow{d} Q x$.

Premultiply (4.19) by $\sqrt{T}$ and rearrange as

$$
\begin{equation*}
\sqrt{T}\left(\hat{\beta}_{L S}-\beta_{0}\right)=\left(\frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right)^{-1} \frac{\sqrt{T}}{T} \sum_{t=1}^{T} x_{t} u_{t} \tag{4.28}
\end{equation*}
$$

If the first term on the right hand side converges in probability to a finite matrix (as assumed in (4.23)), and the vector of random variables $x_{t} u_{t}$ satisfies a central limit theorem, then

$$
\begin{gather*}
\sqrt{T}\left(\hat{\beta}_{L S}-\beta_{0}\right) \xrightarrow{d} N\left(0, \Sigma_{x x}^{-1} S_{0} \Sigma_{x x}^{-1}\right), \text { where } \\
\Sigma_{x x}=\frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime} \text { and } S_{0}=\operatorname{Cov}\left(\frac{\sqrt{T}}{T} \sum_{t=1}^{T} x_{t} u_{t}\right) .
\end{gather*}
$$

The last matrix in the covariance matrix does not need to be transposed since it is symmetric (since $\Sigma_{x x}$ is). This general expression is valid for both autocorrelated and heteroskedastic residuals-all such features are loaded into the $S_{0}$ matrix. Note that $S_{0}$ is the variance-covariance matrix of $\sqrt{T}$ times a sample average (of the vector of random variables $x_{t} u_{t}$ ), which can be complicated to specify and to estimate. In simple cases, we can derive what it is. To do so, we typically need to understand the properties of the residuals. Are they autocorrelated and/or heteroskedastic? In other cases we will have to use some kind of "non-parametric" approach to estimate it.

A common approach is to estimate $\Sigma_{x x}$ by $\Sigma_{t=1}^{T} x_{t} x_{t}^{\prime} / T$ and use the Newey-West estimator of $S_{0}$

### 4.5.1 Special Case: Classical LS assumptions

Reference: Greene (2000) 9.4 or Hamilton (1994) 8.2.
We can recover the classical expression for the covariance, $\sigma^{2} \Sigma_{x x}^{-1}$, if we assume that the regressors are stochastic, but require that $x_{t}$ is independent of all $u_{t+s}$ and that $u_{t}$ is iid. It rules out, for instance, that $u_{t}$ and $x_{t-2}$ are correlated and also that the variance of $u_{t}$ depends on $x_{t}$. Expand the expression for $S_{0}$ as Expand the expression for $S_{0}$ as

$$
\begin{align*}
S_{0} & =\mathrm{E}\left(\frac{\sqrt{T}}{T} \sum_{t=1}^{T} x_{t} u_{t}\right)\left(\frac{\sqrt{T}}{T} \sum_{t=1}^{T} u_{t} x_{t}^{\prime}\right)  \tag{4.30}\\
& =\frac{1}{T} \mathrm{E}\left(\ldots+x_{s-1} u_{s-1}+x_{s} u_{s}+\ldots\right)\left(\ldots+u_{s-1} x_{s-1}^{\prime}+u_{s} x_{s}^{\prime}+\ldots\right) .
\end{align*}
$$

Note that
$\mathrm{E} x_{t-s} u_{t-s} u_{t} x_{t}^{\prime}=\mathrm{E} x_{t-s} x_{t}^{\prime} \mathrm{E} u_{t-s} u_{t}$ (since $u_{t}$ and $x_{t-s}$ independent)

$$
=\left\{\begin{array}{l}
0 \text { if } s \neq 0\left(\text { since } \mathrm{E} u_{t-s} u_{t}=0 \text { by iid } u_{t}\right)  \tag{4.31}\\
\mathrm{E} x_{t} x_{t}^{\prime} \mathrm{E} u_{t} u_{t} \text { else. }
\end{array}\right.
$$

This means that all cross terms (involving different observations) drop out and that we
can write

$$
\begin{align*}
S_{0} & =\frac{1}{T} \sum_{t=1}^{T} \mathrm{E} x_{t} x_{t}^{\prime} \mathrm{E} u_{t}^{2}  \tag{4.32}\\
& =\sigma^{2} \frac{1}{T} \mathrm{E} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\left(\text { since } u_{t} \text { is iid and } \sigma^{2}=\mathrm{E} u_{t}^{2}\right)  \tag{4.33}\\
& =\sigma^{2} \Sigma_{x x} \tag{4.34}
\end{align*}
$$

Using this in (4.29) gives

$$
\text { Asymptotic } \operatorname{Cov}\left[\sqrt{T}\left(\hat{\beta}_{L S}-\beta_{0}\right)\right]=\Sigma_{x x}^{-1} S_{0} \Sigma_{x x}^{-1}=\Sigma_{x x}^{-1} \sigma^{2} \Sigma_{x x} \Sigma_{x x}^{-1}=\sigma^{2} \Sigma_{x x}^{-1}
$$

### 4.5.2 Special Case: White's Heteroskedasticity

Reference: Greene (2000) 12.2 and Davidson and MacKinnon (1993) 16.2.
This section shows that the classical LS formula for the covariance matrix is valid even if the errors are heteroskedastic-provided the heteroskedasticity is independent of the regressors.

The only difference compared with the classical LS assumptions is that $u_{t}$ is now allowed to be heteroskedastic, but this heteroskedasticity is not allowed to depend on the moments of $x_{t}$. This means that (4.32) holds, but (4.33) does not since $\mathrm{E} u_{t}^{2}$ is not the same for all $t$.

However, we can still simplify (4.32) a bit more. We assumed that $\mathrm{E} x_{t} x_{t}^{\prime}$ and $\mathrm{E} u_{t}^{2}$ (which can both be time varying) are not related to each other, so we could perhaps multiply $\mathrm{E} x_{t} x_{t}^{\prime}$ by $\Sigma_{t=1}^{T} \mathrm{E} u_{t}^{2} / T$ instead of by $\mathrm{E} u_{t}^{2}$. This is indeed true asymptotically-where any possible "small sample" relation between $\mathrm{E} x_{t} x_{t}^{\prime}$ and $\mathrm{E} u_{t}^{2}$ must wash out due to the assumptions of independence (which are about population moments).

In large samples we therefore have

$$
\begin{align*}
S_{0} & =\left(\frac{1}{T} \sum_{t=1}^{T} \mathrm{E} u_{t}^{2}\right)\left(\frac{1}{T} \sum_{t=1}^{T} \mathrm{E} x_{t} x_{t}^{\prime}\right) \\
& =\left(\frac{1}{T} \sum_{t=1}^{T} \mathrm{E} u_{t}^{2}\right)\left(\mathrm{E} \frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right) \\
& =\omega^{2} \Sigma_{x x} \tag{4.35}
\end{align*}
$$

where $\omega^{2}$ is a scalar. This is very similar to the classical LS case, except that $\omega^{2}$ is the average variance of the residual rather than the constant variance. In practice, the estimator of $\omega^{2}$ is the same as the estimator of $\sigma^{2}$, so we can actually apply the standard LS formulas in this case.

This is the motivation for why White's test for heteroskedasticity makes sense: if the heteroskedasticity is not correlated with the regressors, then the standard LS formula is correct (provided there is no autocorrelation).

### 4.6 Inference

Consider some estimator, $\hat{\beta}_{k \times 1}$, with an asymptotic normal distribution

$$
\begin{equation*}
\sqrt{T}\left(\hat{\beta}-\beta_{0}\right) \xrightarrow{d} N(0, V) . \tag{4.36}
\end{equation*}
$$

Suppose we want to test the null hypothesis that the $s$ linear restrictions $R \beta_{0}=r$ hold, where $R$ is an $s \times k$ matrix and $r$ is an $s \times 1$ vector. If the null hypothesis is true, then

$$
\begin{equation*}
\sqrt{T}(R \hat{\beta}-r) \xrightarrow{d} N\left(0, R V R^{\prime}\right) \tag{4.37}
\end{equation*}
$$

since the $s$ linear combinations are linear combinations of random variables with an asymptotic normal distribution as in (4.37).

Remark 11 If the $n \times 1$ vector $x \sim N(0, \Sigma)$, then $x^{\prime} \Sigma^{-1} x \sim \chi_{n}^{2}$.

Remark 12 From the previous remark and Remark (9), it follows that if the $n \times 1$ vector $x \xrightarrow{d} N(0, \Sigma)$, then $x^{\prime} \Sigma^{-1} x \xrightarrow{d} \chi_{n}^{2}$.

From this remark, it follows that if the null hypothesis, $R \beta_{0}=r$, is true, then Wald test statistics converges in distribution to a $\chi_{s}^{2}$ variable

$$
\begin{equation*}
T(R \hat{\beta}-r)^{\prime}\left(R V R^{\prime}\right)^{-1}(R \hat{\beta}-r) \xrightarrow{d} \chi_{s}^{2} \tag{4.38}
\end{equation*}
$$

Values of the test statistics above the $x \%$ critical value of the $\chi_{s}^{2}$ distribution mean that we reject the null hypothesis at the $x \%$ significance level.

When there is only one restriction $(s=1)$, then $\sqrt{T}(R \hat{\beta}-r)$ is a scalar, so the test can equally well be based on the fact that

$$
\frac{\sqrt{T}(R \hat{\beta}-r)}{R V R^{\prime}} \xrightarrow{d} N(0,1)
$$

In this case, we should reject the null hypothesis if the test statistics is either very low (negative) or very high (positive). In particular, let $\Phi()$ be the standard normal cumulative distribution function. We then reject the null hypothesis at the $x \%$ significance level if the test statistics is below $x_{L}$ such that $\Phi\left(x_{L}\right)=(x / 2) \%$ or above $x_{H}$ such that $\Phi\left(x_{H}\right)=$ $1-(x / 2) \%$ (that is with $(x / 2) \%$ of the probability mass in each tail).

Example $13\left(T R^{2} /\left(1-R^{2}\right)\right.$ as a test of the regression.) Recall from (4.15)-(4.16) that $R^{2}=\widehat{\operatorname{Var}}\left(\hat{y}_{t}\right) / \widehat{\operatorname{Var}}\left(y_{t}\right)=1-\widehat{\operatorname{Var}}\left(\hat{u}_{t}\right) / \widehat{\operatorname{Var}}\left(y_{t}\right)$, where $\hat{y}_{t}$ and $\hat{u}_{t}$ are the fitted value and residual respectively. We therefore get

$$
T R^{2} /\left(1-R^{2}\right)=T \widehat{\operatorname{Var}}\left(\hat{y}_{t}\right) / \widehat{\operatorname{Var}}\left(\hat{u}_{t}\right) .
$$

To simplify the algebra, assume that both $y_{t}$ and $x_{t}$ are demeaned and that no intercept is used. (We get the same results, but after more work, if we relax this assumption.) In this case, $\hat{y}_{t}=x_{t}^{\prime} \hat{\beta}$, so we can rewrite the previous eqiuation as

$$
T R^{2} /\left(1-R^{2}\right)=T \hat{\beta}^{\prime} \Sigma_{x x} \hat{\beta}^{\prime} / \widehat{\operatorname{Var}}\left(\hat{u}_{t}\right)
$$

This is identical to (4.38) when $R=I_{k}$ and $r=\mathbf{0}_{k \times 1}$ and the classical LS assumptions are fulfilled (so $V=\operatorname{Var}\left(\hat{u}_{t}\right) \Sigma_{x x}^{-1}$ ). The $T R^{2} /\left(1-R^{2}\right)$ is therefore a $\chi_{k}^{2}$ distributed statistics for testing if all the slope coefficients are zero.

Example 14 ( $F$ version of the test.) There is also an $F_{k, T-k}$ version of the test in the previous example: $\left[R^{2} / k\right] /\left[\left(1-R^{2}\right) /(T-k)\right]$. Note that $k$ times an $F_{k, T-k}$ variable converges to a $\chi_{k}^{2}$ variable as $T-k \rightarrow \infty$. This means that the $\chi_{k}^{2}$ form in the previous example can be seen as an asymptotic version of the (more common) F form.

### 4.6.1 On F Tests*

F tests are sometimes used instead of chi-square tests. However, F tests rely on very special assumptions and typically converge to chi-square tests as the sample size increases.

There are therefore few compelling theoretical reasons for why we should use F tests. ${ }^{2}$ This section demonstrates that point.

Remark 15 If $Y_{1} \sim \chi_{n_{1}}^{2}, Y_{2} \sim \chi_{n_{2}}^{2}$, and if $Y_{1}$ and $Y_{2}$ are independent, then $Z=$ $\left(Y_{1} / n_{1}\right) /\left(Y_{1} / n_{1}\right) \sim F_{n_{1}, n_{2}}$. As $n_{2} \rightarrow \infty, n_{1} Z \xrightarrow{d} \chi_{n_{1}}^{2}$ (essentially because the denominator in $Z$ is then equal to its expected value).

To use the F test to test $s$ linear restrictions $R \beta_{0}=r$, we need to assume that the small sample distribution of the estimator is normal, $\sqrt{T}\left(\hat{\beta}-\beta_{0}\right) \sim N\left(0, \sigma^{2} W\right)$, where $\sigma^{2}$ is a scalar and $W$ a known matrix. This would follow from an assumption that the residuals are normally distributed and that we either consider the distribution conditional on the regressors or that the regressors are deterministic. In this case $W=\Sigma_{x x}^{-1}$.

Consider the test statistics

$$
F=T(R \hat{\beta}-r)^{\prime}\left(R \hat{\sigma}^{2} W R^{\prime}\right)^{-1}(R \hat{\beta}-r) / s
$$

This is similar to (4.38), expect that we use the estimated covariance matrix $\hat{\sigma}^{2} W$ instead of the true $\sigma^{2} W$ (recall, $W$ is assumed to be known) and that we have divided by the number of restrictions, $s$. Multiply and divide this expressions by $\sigma^{2}$

$$
F=\frac{T(R \hat{\beta}-r)^{\prime}\left(R \sigma^{2} W R^{\prime}\right)^{-1}(R \hat{\beta}-r) / s}{\hat{\sigma}^{2} / \sigma^{2}}
$$

The numerator is an $\chi_{s}^{2}$ variable divided by its degrees of freedom, $s$. The denominator can be written $\hat{\sigma}^{2} / \sigma^{2}=\Sigma\left(\hat{u}_{t} / \sigma\right)^{2} / T$, where $\hat{u}_{t}$ are the fitted residuals. Since we just assumed that $u_{t}$ are iid $N\left(0, \sigma^{2}\right)$, the denominator is an $\chi_{T}^{2}$ variable divided by its degrees of freedom, $T$. It can also be shown that the numerator and denominator are independent (essentially because the fitted residuals are orthogonal to the regressors), so $F$ is an $F_{s, T}$ variable.

We need indeed very strong assumptions to justify the $F$ distributions. Moreover, as $T \rightarrow \infty, s F \xrightarrow{d} \chi_{n}^{2}$ which is the Wald test-which do not need all these assumptions.

[^2]
### 4.7 Diagnostic Tests of Autocorrelation, Heteroskedasticity, and Normality*

Reference: Greene (2000) 12.3, 13.5 and 9.7; Johnston and DiNardo (1997) 6; and Pindyck and Rubinfeld (1997) 6, Patterson (2000) 5

LS and IV are still consistent even if the residuals are autocorrelated, heteroskedastic, and/or non-normal, but the traditional expression for the variance of the parameter estimators is invalid. It is therefore important to investigate the properties of the residuals.

We would like to test the properties of the true residuals, $u_{t}$, but these are unobservable. We can instead use residuals from a consistent estimator as approximations, since the approximation error then goes to zero as the sample size increases. The residuals from an estimator are

$$
\begin{align*}
\hat{u}_{t} & =y_{t}-x_{t}^{\prime} \hat{\beta} \\
& =x_{t}^{\prime}\left(\beta_{0}-\hat{\beta}\right)+u_{t} . \tag{4.39}
\end{align*}
$$

If $\operatorname{plim} \hat{\beta}=\beta_{0}$, then $\hat{u}_{t}$ converges in probability to the true residual ("pointwise consistency"). It therefore makes sense to use $\hat{u}_{t}$ to study the (approximate) properties of $u_{t}$. We want to understand if $u_{t}$ are autocorrelated and/or heteroskedastic, since this affects the covariance matrix of the least squares estimator and also to what extent least squares is efficient. We might also be interested in studying if the residuals are normally distributed, since this also affects the efficiency of least squares (remember that LS is MLE is the residuals are normally distributed).

It is important that the fitted residuals used in the diagnostic tests are consistent. With poorly estimated residuals, we can easily find autocorrelation, heteroskedasticity, or nonnormality even if the true residuals have none of these features

### 4.7.1 Autocorrelation

Let $\hat{\rho}_{s}$ be the estimate of the $s$ th autocorrelation coefficient of some variable, for instance, the fitted residuals. The sampling properties of $\hat{\rho}_{s}$ are complicated, but there are several useful large sample results for Gaussian processes (these results typically carry over to processes which are similar to the Gaussian-a homoskedastic process with finite 6th moment is typically enough). When the true autocorrelations are all zero (not $\rho_{0}$, of
course), then for any $i$ and $j$ different from zero

$$
\sqrt{T}\left[\begin{array}{l}
\hat{\rho}_{i}  \tag{4.40}\\
\hat{\rho}_{j}
\end{array}\right] \rightarrow^{d} N\left(\left[\begin{array}{l}
0 \\
0
\end{array}\right],\left[\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right]\right) .
$$

This result can be used to construct tests for both single autocorrelations (t-test or $\chi^{2}$ test) and several autocorrelations at once ( $\chi^{2}$ test)

Example 16 (t-test) We want to test the hypothesis that $\rho_{1}=0$. Since the $N(0,1)$ distribution has 5\% of the probability mass below -1.65 and another 5\% above 1.65, we can reject the null hypothesis at the $10 \%$ level if $\sqrt{T}\left|\hat{\rho}_{1}\right|>1.65$. With $T=100$, we therefore need $\left|\hat{\rho}_{1}\right|>1.65 / \sqrt{100}=0.165$ for rejection, and with $T=1000$ we need $\left|\hat{\rho}_{1}\right|>1.65 / \sqrt{1000} \approx 0.0 .53$.

The Box-Pierce test follows directly from the result in (4.40), since it shows that $\sqrt{T} \hat{\rho}_{i}$ and $\sqrt{T} \hat{\rho}_{j}$ are iid $\mathrm{N}(0,1)$ variables. Therefore, the sum of the square of them is distributed as an $\chi^{2}$ variable. The test statistics typically used is

$$
\begin{equation*}
Q_{L}=T \sum_{s=1}^{L} \hat{\rho}_{s}^{2} \rightarrow^{d} \chi_{L}^{2} \tag{4.41}
\end{equation*}
$$

Example 17 (Box-Pierce) Let $\hat{\rho}_{1}=0.165$, and $T=100$, so $Q_{1}=100 \times 0.165^{2}=$ 2.72. The $10 \%$ critical value of the $\chi_{1}^{2}$ distribution is 2.71 , so the null hypothesis of no autocorrelation is rejected.

The choice of lag order in (4.41), $L$, should be guided by theoretical considerations, but it may also be wise to try different values. There is clearly a trade off: too few lags may miss a significant high-order autocorrelation, but too many lags can destroy the power of the test (as the test statistics is not affected much by increasing $L$, but the critical values increase).

Example 18 (Residuals follow an $A R(1)$ process) If $u_{t}=0.9 u_{t-1}+\varepsilon_{t}$, then the true autocorrelation coefficients are $\rho_{j}=0.9^{j}$.

A common test of the serial correlation of residuals from a regression is the DurbinWatson test

$$
\begin{equation*}
d=2\left(1-\hat{\rho}_{1}\right) \tag{4.42}
\end{equation*}
$$

where the null hypothesis of no autocorrelation is

$$
\begin{aligned}
& \text { not rejected if } d>d_{\text {upper }}^{*} \\
& \text { rejected if } d<d_{\text {lower }}^{*} \text { (in favor of positive autocorrelation) } \\
& \text { else inconclusive }
\end{aligned}
$$

where the upper and lower critical values can be found in tables. (Use $4-d$ to let negative autocorrelation be the alternative hypothesis.) This test is typically not useful when lagged dependent variables enter the right hand side ( $d$ is biased towards showing no autocorrelation). Note that DW tests only for first-order autocorrelation.

Example 19 (Durbin-Watson.) With $\hat{\rho}_{1}=0.2$ we get $d=1.6$. For large samples, the $5 \%$ critical value is $d_{\text {lower }}^{*} \approx 1.6$, so $\hat{\rho}_{1}>0.2$ is typically considered as evidence of positive autocorrelation.

The fitted residuals used in the autocorrelation tests must be consistent in order to interpret the result in terms of the properties of the true residuals. For instance, an excluded autocorrelated variable will probably give autocorrelated fitted residuals-and also make the coefficient estimator inconsistent (unless the excluded variable is uncorrelated with the regressors). Only when we know that the model is correctly specified can we interpret a finding of autocorrelated residuals as an indication of the properties of the true residuals.

### 4.7.2 Heteroskedasticity

Remark 20 (Kronecker product.) If $A$ and $B$ are matrices, then

$$
A \otimes B=\left[\begin{array}{ccc}
a_{11} B & \cdots & a_{1 n} B \\
\vdots & & \vdots \\
a_{m 1} B & \cdots & a_{m n} B
\end{array}\right]
$$

Example 21 Let $x_{1}$ and $x_{2}$ be scalars. Then

White's test for heteroskedasticity tests the null hypothesis of homoskedasticity against the kind of heteroskedasticity which can be explained by the levels, squares, and cross products of the regressors. Let $w_{t}$ be the unique elements in $x_{t} \otimes x_{t}$, where we have added a constant to $x_{t}$ if there was not one from the start. Run a regression of the squared fitted LS residuals on $w_{t}$

$$
\begin{equation*}
\hat{u}_{t}^{2}=w_{t}^{\prime} \gamma+\varepsilon_{t} \tag{4.43}
\end{equation*}
$$

and test if all elements (except the constant) in $\gamma$ are zero (with a $\chi^{2}$ or $F$ test). The reason for this specification is that if $u_{t}^{2}$ is uncorrelated with $x_{t} \otimes x_{t}$, then the usual LS covariance matrix applies.

Breusch-Pagan's test is very similar, except that the vector $w_{t}$ in (4.43) can be any vector which is thought of as useful for explaining the heteroskedasticity. The null hypothesis is that the variance is constant, which is tested against the alternative that the variance is some function of $w_{t}$.

The fitted residuals used in the heteroskedasticity tests must be consistent in order to interpret the result in terms of the properties of the true residuals. For instance, if some of the of elements in $w_{t}$ belong to the regression equation, but are excluded, then fitted residuals will probably fail these tests.

### 4.7.3 Normality

We often make the assumption of normally distributed errors, for instance, in maximum likelihood estimation. This assumption can be tested by using the fitted errors. This works since moments estimated from the fitted errors are consistent estimators of the moments of the true errors. Define the degree of skewness and excess kurtosis for a variable $z_{t}$ (could be the fitted residuals) as

$$
\begin{align*}
& \hat{\theta}_{3}=\frac{1}{T} \sum_{t=1}^{T}\left(z_{t}-\bar{z}\right)^{3} / \hat{\sigma}^{3}  \tag{4.44}\\
& \hat{\theta}_{4}=\frac{1}{T} \sum_{t=1}^{T}\left(z_{t}-\bar{z}\right)^{4} / \hat{\sigma}^{4}-3 \tag{4.45}
\end{align*}
$$

where $\bar{z}$ is the sample mean and $\hat{\sigma}^{2}$ is the estimated variance.
Remark $22\left(\chi^{2}(n)\right.$ distribution.) If $x_{i}$ are independent $N\left(0, \sigma_{i}^{2}\right)$ variables, then $\Sigma_{i=1}^{n} x_{i}^{2} / \sigma_{i}^{2} \sim$


Figure 4.1: This figure shows a histogram from 100 draws of iid uniformly [0,1] distributed variables.

## $\chi^{2}(n)$.

In a normal distribution, the true values are zero and the test statistics $\hat{\theta}_{3}$ and $\hat{\theta}_{4}$ are themselves normally distributed with zero covariance and variances $6 / T$ and $24 / T$, respectively (straightforward, but tedious, to show). Therefore, under the null hypothesis of a normal distribution, $T \hat{\theta}_{3}^{2} / 6$ and $T \hat{\theta}_{4}^{2} / 24$ are independent and both asymptotically distributed as $\chi^{2}(1)$, so the sum is asymptotically a $\chi^{2}(2)$ variable

$$
\begin{equation*}
W=T\left(\hat{\theta}_{3}^{2} / 6+\hat{\theta}_{4}^{2} / 24\right) \rightarrow^{d} \chi^{2}(2) \tag{4.46}
\end{equation*}
$$

This is the Jarque and Bera test of normality.

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## 5 Instrumental Variable Method

Reference: Greene (2000) 9.5 and 16.1-2
Additional references: Hayashi (2000) 3.1-4; Verbeek (2000) 5.1-4; Hamilton (1994) 8.2; and Pindyck and Rubinfeld (1997) 7

### 5.1 Consistency of Least Squares or Not?

Consider the linear model

$$
\begin{equation*}
y_{t}=x_{t}^{\prime} \beta_{0}+u_{t} \tag{5.1}
\end{equation*}
$$

where $y_{t}$ and $u_{t}$ are scalars, $x_{t}$ a $k \times 1$ vector, and $\beta_{0}$ is a $k \times 1$ vector of the true coefficients. The least squares estimator is

$$
\begin{align*}
\hat{\beta}_{L S} & =\left(\frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} x_{t} y_{t}  \tag{5.2}\\
& =\beta_{0}+\left(\frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} x_{t} u_{t} \tag{5.3}
\end{align*}
$$

where we have used (5.1) to substitute for $y_{t}$. The probability limit is

$$
\begin{equation*}
\operatorname{plim} \hat{\beta}_{L S}-\beta_{0}=\left(\operatorname{plim} \frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right)^{-1} \operatorname{plim} \frac{1}{T} \sum_{t=1}^{T} x_{t} u_{t} \tag{5.4}
\end{equation*}
$$

In many cases the law of large numbers applies to both terms on the right hand side. The first term is typically a matrix with finite elements and the second term is the covariance of the regressors and the true residuals. This covariance must be zero for LS to be consistent.

### 5.2 Reason 1 for IV: Measurement Errors

Reference: Greene (2000) 9.5.

Suppose the true model is

$$
\begin{equation*}
y_{t}^{*}=x_{t}^{* \prime} \beta_{0}+u_{t}^{*} \tag{5.5}
\end{equation*}
$$

Data on $y_{t}^{*}$ and $x_{t}^{*}$ is not directly observable, so we instead run the regression

$$
\begin{equation*}
y_{t}=x_{t}^{\prime} \beta+u_{t} \tag{5.6}
\end{equation*}
$$

where $y_{t}$ and $x_{t}$ are proxies for the correct variables (the ones that the model is true for). We can think of the difference as measurement errors

$$
\begin{align*}
& y_{t}=y_{t}^{*}+v_{t}^{y} \text { and }  \tag{5.7}\\
& x_{t}=x_{t}^{*}+v_{t}^{x}, \tag{5.8}
\end{align*}
$$

where the errors are uncorrelated with the true values and the "true" residual $u_{t}^{*}$.
Use (5.7) and (5.8) in (5.5)

$$
\begin{align*}
y_{t}-v_{t}^{y} & =\left(x_{t}-v_{t}^{x}\right)^{\prime} \beta_{0}+u_{t}^{*} \text { or } \\
y_{t} & =x_{t}^{\prime} \beta_{0}+\varepsilon_{t} \text { where } \varepsilon_{t}=-v_{t}^{x \prime} \beta_{0}+v_{t}^{y}+u_{t}^{*} . \tag{5.9}
\end{align*}
$$

Suppose that $x_{t}^{*}$ is a measured with error. From (5.8) we see that $v_{t}^{x}$ and $x_{t}$ are correlated, so LS on (5.9) is inconsistent in this case. To make things even worse, measurement errors in only one of the variables typically affect all the coefficient estimates.

To illustrate the effect of the error, consider the case when $x_{t}$ is a scalar. Then, the probability limit of the LS estimator of $\beta$ in (5.9) is

$$
\begin{align*}
\operatorname{plim} \hat{\beta}_{L S} & =\operatorname{Cov}\left(y_{t}, x_{t}\right) / \operatorname{Var}\left(x_{t}\right) \\
& =\operatorname{Cov}\left(x_{t}^{*} \beta_{0}+u_{t}^{*}, x_{t}\right) / \operatorname{Var}\left(x_{t}\right) \\
& =\operatorname{Cov}\left(x_{t} \beta_{0}-v_{t}^{x} \beta_{0}+u_{t}^{*}, x_{t}\right) / \operatorname{Var}\left(x_{t}\right) \\
& =\frac{\operatorname{Cov}\left(x_{t} \beta_{0}, x_{t}\right)+\operatorname{Cov}\left(-v_{t}^{x} \beta_{0}, x_{t}\right)+\operatorname{Cov}\left(u_{t}^{*}, x_{t}\right)}{\operatorname{Var}\left(x_{t}\right)} \\
& =\frac{\operatorname{Var}\left(x_{t}\right)}{\operatorname{Var}\left(x_{t}\right)} \beta_{0}+\frac{\operatorname{Cov}\left(-v_{t}^{x} \beta_{0}, x_{t}^{*}-v_{t}^{x}\right)}{\operatorname{Var}\left(x_{t}\right)} \\
& =\beta_{0}-\beta_{0} \operatorname{Var}\left(v_{t}^{x}\right) / \operatorname{Var}\left(x_{t}\right) \\
& =\beta_{0}\left[1-\frac{\operatorname{Var}\left(v_{t}^{x}\right)}{\operatorname{Var}\left(x_{t}^{*}\right)+\operatorname{Var}\left(v_{t}^{x}\right)}\right] . \tag{5.10}
\end{align*}
$$

since $x_{t}^{*}$ and $v_{t}^{x}$ are uncorrelated with $u_{t}^{*}$ and with each other. This shows that $\hat{\beta}_{L S}$ goes to zero as the measurement error becomes relatively more volatile compared with the true value. This makes a lot of sense, since when the measurement error is very large then the regressor $x_{t}$ is dominated by noise that has nothing to do with the dependent variable.

Suppose instead that only $y_{t}^{*}$ is measured with error. This not a big problem since this measurement error is uncorrelated with the regressor, so the consistency of least squares is not affected. In fact, a measurement error in the dependent variable is like increasing the variance in the residual.

### 5.3 Reason 2 for IV: Simultaneous Equations Bias (and Inconsistency)

Suppose economic theory tells you that the structural form of the $m$ endogenous variables, $y_{t}$, and the $k$ predetermined (exogenous) variables, $z_{t}$, is

$$
\begin{equation*}
F y_{t}+G z_{t}=u_{t}, \text { where } u_{t} \text { is iid with } \mathrm{E} u_{t}=\mathbf{0} \text { and } \operatorname{Cov}\left(u_{t}\right)=\Sigma, \tag{5.11}
\end{equation*}
$$

where $F$ is $m \times m$, and $G$ is $m \times k$. The disturbances are assumed to be uncorrelated with the predetermined variables, $\mathrm{E}\left(z_{t} u_{t}^{\prime}\right)=\mathbf{0}$.

Suppose $F$ is invertible. Solve for $y_{t}$ to get the reduced form

$$
\begin{align*}
y_{t} & =-F^{-1} G z_{t}+F^{-1} u_{t}  \tag{5.12}\\
& =\Pi z_{t}+\varepsilon_{t}, \text { with } \operatorname{Cov}\left(\varepsilon_{t}\right)=\Omega . \tag{5.13}
\end{align*}
$$

The reduced form coefficients, $П$, can be consistently estimated by LS on each equation since the exogenous variables $z_{t}$ are uncorrelated with the reduced form residuals (which are linear combinations of the structural residuals). The fitted residuals can then be used to get an estimate of the reduced form covariance matrix.

The $j$ th line of the structural form (5.11) can be written

$$
\begin{equation*}
F_{j} y_{t}+G_{j} z_{t}=u_{j t} \tag{5.14}
\end{equation*}
$$

where $F_{j}$ and $G_{j}$ are the $j$ th rows of $F$ and $G$, respectively. Suppose the model is normalized so that the coefficient on $y_{j t}$ is one (otherwise, divide (5.14) with this coefficient).

Then, rewrite (5.14) as

$$
\begin{align*}
y_{j t} & =-G_{j 1} \tilde{z}_{t}-F_{j 1} \tilde{y}_{t}+u_{j t} \\
& =x_{t}^{\prime} \beta+u_{j t}, \text { where } x_{t}^{\prime}=\left[\tilde{z}_{t}^{\prime}, \tilde{y}_{t}^{\prime}\right], \tag{5.15}
\end{align*}
$$

where $\tilde{z}_{t}$ and $\tilde{y}_{t}$ are the exogenous and endogenous variables that enter the $j$ th equation, which we collect in the $x_{t}$ vector to highlight that (5.15) looks like any other linear regression equation. The problem with (5.15), however, is that the residual is likely to be correlated with the regressors, so the LS estimator is inconsistent. The reason is that a shock to $u_{j t}$ influences $y_{j t}$, which in turn will affect some other endogenous variables in the system (5.11). If any of these endogenous variable are in $x_{t}$ in (5.15), then there is a correlation between the residual and (some of) the regressors

Note that the concept of endogeneity discussed here only refers to contemporaneous endogeneity as captured by off-diagonal elements in $F$ in (5.11). The vector of predetermined variables, $z_{t}$, could very well include lags of $y_{t}$ without affecting the econometric endogeneity problem.

Example 1 (Supply and Demand. Reference: GR 16, Hamilton 9.1.) Consider the simplest simultaneous equations model for supply and demand on a market. Supply is

$$
q_{t}=\gamma p_{t}+u_{t}^{s}, \gamma>0,
$$

and demand is

$$
q_{t}=\beta p_{t}+\alpha A_{t}+u_{t}^{d}, \beta<0
$$

where $A_{t}$ is an observable demand shock (perhaps income). The structural form is therefore

$$
\left[\begin{array}{ll}
1 & -\gamma \\
1 & -\beta
\end{array}\right]\left[\begin{array}{l}
q_{t} \\
p_{t}
\end{array}\right]+\left[\begin{array}{l}
0 \\
-\alpha
\end{array}\right] A_{t}=\left[\begin{array}{l}
u_{t}^{s} \\
u_{t}^{d}
\end{array}\right]
$$

The reduced form is

$$
\left[\begin{array}{l}
q_{t} \\
p_{t}
\end{array}\right]=\left[\begin{array}{l}
\pi_{11} \\
\pi_{21}
\end{array}\right] A_{t}+\left[\begin{array}{l}
\varepsilon_{1 t} \\
\varepsilon_{2 t}
\end{array}\right] .
$$

If we knew the structural form, then we can solve for $q_{t}$ and $p_{t}$ to get the reduced form in
terms of the structural parameters

$$
\left[\begin{array}{c}
q_{t} \\
p_{t}
\end{array}\right]=\left[\begin{array}{c}
-\frac{\gamma}{\beta-\gamma} \alpha \\
-\frac{1}{\beta-\gamma} \alpha
\end{array}\right] A_{t}+\left[\begin{array}{cc}
\frac{\beta}{\beta-\gamma} & -\frac{\gamma}{\beta-\gamma} \\
\frac{1}{\beta-\gamma} & -\frac{1}{\beta-\gamma}
\end{array}\right]\left[\begin{array}{c}
u_{t}^{s} \\
u_{t}^{d}
\end{array}\right] .
$$

Example 2 (Supply equation with LS.) Suppose we try to estimate the supply equation in Example 1 by LS, that is, we run the regression

$$
q_{t}=\theta p_{t}+\varepsilon_{t}
$$

If data is generated by the model in Example 1, then the reduced form shows that $p_{t}$ is correlated with $u_{t}^{s}$, so we cannot hope that $L S$ will be consistent. In fact, when both $q_{t}$ and $p_{t}$ have zero means, then the probability limit of the LS estimator is

$$
\begin{aligned}
\operatorname{plim} \hat{\theta} & =\frac{\operatorname{Cov}\left(q_{t}, p_{t}\right)}{\operatorname{Var}\left(p_{t}\right)} \\
& =\frac{\operatorname{Cov}\left(\frac{\gamma \alpha}{\gamma-\beta} A_{t}+\frac{\gamma}{\gamma-\beta} u_{t}^{d}-\frac{\beta}{\gamma-\beta} u_{t}^{s}, \frac{\alpha}{\gamma-\beta} A_{t}+\frac{1}{\gamma-\beta} u_{t}^{d}-\frac{1}{\gamma-\beta} u_{t}^{d}\right)}{\operatorname{Var}\left(\frac{\alpha}{\gamma-\beta} A_{t}+\frac{1}{\gamma-\beta} u_{t}^{d}-\frac{1}{\gamma-\beta} u_{t}^{s}\right),}
\end{aligned}
$$

where the second line follows from the reduced form. Suppose the supply and demand shocks are uncorrelated. In that case we get

$$
\begin{aligned}
\operatorname{plim} \hat{\theta} & =\frac{\frac{\gamma \alpha^{2}}{(\gamma-\beta)^{2}} \operatorname{Var}\left(A_{t}\right)+\frac{\gamma}{(\gamma-\beta)^{2}} \operatorname{Var}\left(u_{t}^{d}\right)+\frac{\beta}{(\gamma-\beta)^{2}} \operatorname{Var}\left(u_{t}^{s}\right)}{\frac{\alpha^{2}}{(\gamma-\beta)^{2}} \operatorname{Var}\left(A_{t}\right)+\frac{1}{(\gamma-\beta)^{2}} \operatorname{Var}\left(u_{t}^{d}\right)+\frac{1}{(\gamma-\beta)^{2}} \operatorname{Var}\left(u_{t}^{s}\right)} \\
& =\frac{\gamma \alpha^{2} \operatorname{Var}\left(A_{t}\right)+\gamma \operatorname{Var}\left(u_{t}^{d}\right)+\beta \operatorname{Var}\left(u_{t}^{s}\right)}{\alpha^{2} \operatorname{Var}\left(A_{t}\right)+\operatorname{Var}\left(u_{t}^{d}\right)+\operatorname{Var}\left(u_{t}^{s}\right)} .
\end{aligned}
$$

First, suppose the supply shocks are zero, $\operatorname{Var}\left(u_{t}^{s}\right)=0$, then $\operatorname{plim} \hat{\theta}=\gamma$, so we indeed estimate the supply elasticity, as we wanted. Think of a fixed supply curve, and a demand curve which moves around. These point of $p_{t}$ and $q_{t}$ should trace out the supply curve. It is clearly $u_{t}^{s}$ that causes a simultaneous equations problem in estimating the supply curve: $u_{t}^{s}$ affects both $q_{t}$ and $p_{t}$ and the latter is the regressor in the supply equation. With no movements in $u_{t}^{s}$ there is no correlation between the shock and the regressor. Second, now suppose instead that the both demand shocks are zero (both $A_{t}=0$ and $\operatorname{Var}\left(u_{t}^{d}\right)=0$ ). Then $\operatorname{plim} \hat{\theta}=\beta$, so the estimated value is not the supply, but the demand elasticity. Not good. This time, think of a fixed demand curve, and a supply curve which moves around.

Example 3 (A flat demand curve.) Suppose we change the demand curve in Example 1 to be infinitely elastic, but to still have demand shocks. For instance, the inverse demand curve could be $p_{t}=\psi A_{t}+u_{t}^{D}$. In this case, the supply and demand is no longer a simultaneous system of equations and both equations could be estimated consistently with LS. In fact, the system is recursive, which is easily seen by writing the system on vector form

$$
\left[\begin{array}{ll}
1 & 0 \\
1 & -\gamma
\end{array}\right]\left[\begin{array}{l}
p_{t} \\
q_{t}
\end{array}\right]+\left[\begin{array}{l}
-\psi \\
0
\end{array}\right] A_{t}=\left[\begin{array}{l}
u_{t}^{D} \\
u_{t}^{s}
\end{array}\right] .
$$

A supply shock, $u_{t}^{s}$, affects the quantity, but this has no affect on the price (the regressor in the supply equation), so there is no correlation between the residual and regressor in the supply equation. A demand shock, $u_{t}^{D}$, affects the price and the quantity, but since quantity is not a regressor in the inverse demand function (only the exogenous $A_{t}$ is) there is no correlation between the residual and the regressor in the inverse demand equation either.

### 5.4 Definition of the IV Estimator-Consistency of IV

Reference: Greene (2000) 9.5; Hamilton (1994) 8.2; and Pindyck and Rubinfeld (1997) 7.

Consider the linear model

$$
\begin{equation*}
y_{t}=x_{t}^{\prime} \beta_{0}+u_{t} \tag{5.16}
\end{equation*}
$$

where $y_{t}$ is a scalar, $x_{t}$ a $k \times 1$ vector, and $\beta_{0}$ is a vector of the true coefficients. If we suspect that $x_{t}$ and $u_{t}$ in (5.16) are correlated, then we may use the instrumental variables (IV) method. To do that, let $z_{t}$ be a $k \times 1$ vector of instruments (as many instruments as regressors; we will later deal with the case when we have more instruments than regressors.) If $x_{t}$ and $u_{t}$ are not correlated, then setting $x_{t}=z_{t}$ gives the least squares (LS) method.

Recall that LS minimizes the variance of the fitted residuals, $\hat{u}_{t}=y_{t}-x_{t}^{\prime} \hat{\beta}_{L S}$. The first order conditions for that optimization problem are

$$
\begin{equation*}
\mathbf{0}_{k x 1}=\frac{1}{T} \sum_{t=1}^{T} x_{t}\left(y_{t}-x_{t}^{\prime} \hat{\beta}_{L S}\right) \tag{5.17}
\end{equation*}
$$

If $x_{t}$ and $u_{t}$ are correlated, then $\operatorname{plim} \hat{\beta}_{L S} \neq \beta_{0}$. The reason is that the probability limit of the right hand side of $(5.17)$ is $\operatorname{Cov}\left(x_{t}, y_{t}-x_{t}^{\prime} \hat{\beta}_{L S}\right)$, which at $\hat{\beta}_{L S}=\beta_{0}$ is non-zero, so the first order conditions (in the limit) cannot be satisfied at the true parameter values. Note that since the LS estimator by construction forces the fitted residuals to be uncorrelated with the regressors, the properties of the LS residuals are of little help in deciding if to use LS or IV.

The idea of the IV method is to replace the first $x_{t}$ in (5.17) with a vector (of similar size) of some instruments, $z_{t}$. The identifying assumption of the IV method is that the instruments are uncorrelated with the residuals (and, as we will see, correlated with the regressors)

$$
\begin{align*}
\mathbf{0}_{k x 1} & =\mathrm{E} z_{t} u_{t}  \tag{5.18}\\
& =\mathrm{E} z_{t}\left(y_{t}-x_{t}^{\prime} \beta_{0}\right) . \tag{5.19}
\end{align*}
$$

The intuition is that the linear model (5.16) is assumed to be correctly specified: the residuals, $u_{t}$, represent factors which we cannot explain, so $z_{t}$ should not contain any information about $u_{t}$.

The sample analogue to (5.19) defines the IV estimator of $\beta$ as ${ }^{1}$

$$
\begin{align*}
& \mathbf{0}_{k x 1}=\frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \hat{\beta}_{I V}\right), \text { or }  \tag{5.20}\\
& \hat{\beta}_{I V}=\left(\frac{1}{T} \sum_{t=1}^{T} z_{t} x_{t}^{\prime}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} z_{t} y_{t} . \tag{5.21}
\end{align*}
$$

It is clearly necessay for $\Sigma z_{t} x_{t}^{\prime} / T$ to have full rank to calculate the IV estimator.
Remark 4 (Probability limit of product) For any random variables $y_{T}$ and $x_{T}$ where plim $y_{T}=a$ and $\operatorname{plim} x_{T}=b$ ( $a$ and $b$ are constants), we have plim $y_{T} x_{T}=a b$.

To see if the IV estimator is consistent, use (5.16) to substitute for $y_{t}$ in (5.20) and take the probability limit

$$
\begin{equation*}
\operatorname{plim} \frac{1}{T} \sum_{t=1}^{T} z_{t} x_{t}^{\prime} \beta_{0}+\operatorname{plim} \frac{1}{T} \sum_{t=1}^{T} z_{t} u_{t}=\operatorname{plim} \frac{1}{T} \sum_{t=1}^{T} z_{t} x_{t}^{\prime} \hat{\beta}_{I V} \tag{5.22}
\end{equation*}
$$

[^3]Two things are required for consistency of the IV estimator, $\operatorname{plim} \hat{\beta}_{I V}=\beta_{0}$. First, that $\operatorname{plim} \Sigma z_{t} u_{t} / T=0$. Provided a law of large numbers apply, this is condition (5.18). Second, that $\operatorname{plim} \Sigma z_{t} x_{t}^{\prime} / T$ has full rank. To see this, suppose $\operatorname{plim} \Sigma z_{t} u_{t} / T=0$ is satisfied. Then, (5.22) can be written

$$
\begin{equation*}
\left(\operatorname{plim} \frac{1}{T} \sum_{t=1}^{T} z_{t} x_{t}^{\prime}\right)\left(\beta_{0}-\operatorname{plim} \hat{\beta}_{I V}\right)=0 . \tag{5.23}
\end{equation*}
$$

If $\operatorname{plim} \Sigma z_{t} x_{t}^{\prime} / T$ has reduced rank, then $\operatorname{plim} \hat{\beta}_{I V}$ does not need to equal $\beta_{0}$ for (5.23) to be satisfied. In practical terms, the first order conditions (5.20) do then not define a unique value of the vector of estimates. If a law of large numbers applies, then plim $\Sigma z_{t} x_{t}^{\prime} / T=$ $\mathrm{E} z_{t} x_{t}^{\prime}$. If both $z_{t}$ and $x_{t}$ contain constants (or at least one of them has zero means), then a reduced rank of $\mathrm{E} z_{t} x_{t}^{\prime}$ would be a consequence of a reduced rank of the covariance matrix of the stochastic elements in $z_{t}$ and $x_{t}$, for instance, that some of the instruments are uncorrelated with all the regressors. This shows that the instruments must indeed be correlated with the regressors for IV to be consistent (and to make sense).

Remark 5 (Second moment matrix) Note that $E z x^{\prime}=E z E x x^{\prime}+\operatorname{Cov}(z, x)$. If $E z=\mathbf{0}$ and/or $E x=\mathbf{0}$, then the second moment matrix is a covariance matrix. Alternatively, suppose both $z$ and $x$ contain constants normalized to unity: $z=\left[1, \tilde{z}^{\prime}\right]^{\prime}$ and $x=\left[1, \tilde{x}^{\prime}\right]^{\prime}$ where $\tilde{z}$ and $\tilde{x}$ are random vectors. We can then write

$$
\begin{aligned}
E z x^{\prime} & =\left[\begin{array}{c}
1 \\
E \tilde{z}
\end{array}\right]\left[\begin{array}{cc}
1 & E \tilde{x}^{\prime}
\end{array}\right]+\left[\begin{array}{cc}
0 & 0 \\
0 & \operatorname{Cov}(\tilde{z}, \tilde{x})
\end{array}\right] \\
& =\left[\begin{array}{cc}
1 & E \tilde{x}^{\prime} \\
E \tilde{z} & E \tilde{z} E \tilde{x}^{\prime}+\operatorname{Cov}(\tilde{z}, \tilde{x})
\end{array}\right] .
\end{aligned}
$$

For simplicity, suppose $\tilde{z}$ and $\tilde{x}$ are scalars. Then Ezx' has reduced rank if $\operatorname{Cov}(\tilde{z}, \tilde{x})=0$, since $\operatorname{Cov}(\tilde{z}, \tilde{x})$ is then the determinant of $E z x^{\prime}$. This is true also when $\tilde{z}$ and $\tilde{x}$ are vectors.

Example 6 (Supply equation with IV.) Suppose we try to estimate the supply equation in Example 1 by IV. The only available instrument is $A_{t}$, so (5.21) becomes

$$
\hat{\gamma}_{I V}=\left(\frac{1}{T} \sum_{t=1}^{T} A_{t} p_{t}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} A_{t} q_{t},
$$

so the probability limit is

$$
\operatorname{plim} \hat{\gamma}_{I V}=\operatorname{Cov}\left(A_{t}, p_{t}\right)^{-1} \operatorname{Cov}\left(A_{t}, q_{t}\right)
$$

since all variables have zero means. From the reduced form in Example 1 we see that

$$
\operatorname{Cov}\left(A_{t}, p_{t}\right)=-\frac{1}{\beta-\gamma} \alpha \operatorname{Var}\left(A_{t}\right) \text { and } \operatorname{Cov}\left(A_{t}, q_{t}\right)=-\frac{\gamma}{\beta-\gamma} \alpha \operatorname{Var}\left(A_{t}\right)
$$

so

$$
\begin{aligned}
\operatorname{pim} \hat{\gamma}_{I V} & =\left[-\frac{1}{\beta-\gamma} \alpha \operatorname{Var}\left(A_{t}\right)\right]^{-1}\left[-\frac{\gamma}{\beta-\gamma} \alpha \operatorname{Var}\left(A_{t}\right)\right] \\
& =\gamma .
\end{aligned}
$$

This shows that $\hat{\gamma}_{I V}$ is consistent.

### 5.4.1 Asymptotic Normality of IV

Little is known about the finite sample distribution of the IV estimator, so we focus on the asymptotic distribution-assuming the IV estimator is consistent.

Remark 7 If $x_{T} \xrightarrow{d} x$ (a random variable) and $\operatorname{plim} Q_{T}=Q$ (a constant matrix), then $Q_{T} x_{T} \xrightarrow{d} Q x$.

Use (5.16) to substitute for $y_{t}$ in (5.20)

$$
\hat{\beta}_{I V}=\beta_{0}+\left(\frac{1}{T} \sum_{t=1}^{T} z_{t} x_{t}^{\prime}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} z_{t} u_{t}
$$

Premultiply by $\sqrt{T}$ and rearrange as

$$
\begin{equation*}
\sqrt{T}\left(\hat{\beta}_{I V}-\beta_{0}\right)=\left(\frac{1}{T} \sum_{t=1}^{T} z_{t} x_{t}^{\prime}\right)^{-1} \frac{\sqrt{T}}{T} \sum_{t=1}^{T} z_{t} u_{t} . \tag{5.25}
\end{equation*}
$$

If the first term on the right hand side converges in probability to a finite matrix (as assumed in in proving consistency), and the vector of random variables $z_{t} u_{t}$ satisfies a
central limit theorem, then

$$
\begin{gather*}
\sqrt{T}\left(\hat{\beta}_{I V}-\beta_{0}\right) \xrightarrow{d} N\left(0, \Sigma_{z x}^{-1} S_{0} \Sigma_{x z}^{-1}\right), \text { where }  \tag{5.26}\\
\Sigma_{z x}=\frac{1}{T} \sum_{t=1}^{T} z_{t} x_{t}^{\prime} \text { and } S_{0}=\operatorname{Cov}\left(\frac{\sqrt{T}}{T} \sum_{t=1}^{T} z_{t} u_{t}\right) .
\end{gather*}
$$

The last matrix in the covariance matrix follows from $\left(\Sigma_{z x}^{-1}\right)^{\prime}=\left(\Sigma_{z x}^{\prime}\right)^{-1}=\Sigma_{x z}^{-1}$. This general expression is valid for both autocorrelated and heteroskedastic residuals-all such features are loaded into the $S_{0}$ matrix. Note that $S_{0}$ is the variance-covariance matrix of $\sqrt{T}$ times a sample average (of the vector of random variables $x_{t} u_{t}$ ).

Example 8 (Choice of instrument in IV, simplest case) Consider the simple regression

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

The asymptotic variance of the IV estimator is

$$
\operatorname{AVar}\left(\sqrt{T}\left(\hat{\beta}_{I V}-\beta_{0}\right)\right)=\operatorname{Var}\left(\frac{\sqrt{T}}{T} \sum_{t=1}^{T} z_{t} u_{t}\right) / \operatorname{Cov}\left(z_{t}, x_{t}\right)^{2}
$$

If $z_{t}$ and $u_{t}$ is serially uncorrelated and independent of each other, then $\operatorname{Var}\left(\Sigma_{t=1}^{T} z_{t} u_{t} / \sqrt{T}\right)=$ $\operatorname{Var}\left(z_{t}\right) \operatorname{Var}\left(u_{t}\right)$. We can then write

$$
\operatorname{AVar}\left(\sqrt{T}\left(\hat{\beta}_{I V}-\beta_{0}\right)\right)=\operatorname{Var}\left(u_{t}\right) \frac{\operatorname{Var}\left(z_{t}\right)}{\operatorname{Cov}\left(z_{t}, x_{t}\right)^{2}}=\frac{\operatorname{Var}\left(u_{t}\right)}{\operatorname{Var}\left(x_{t}\right) \operatorname{Corr}\left(z_{t}, x_{t}\right)^{2}} .
$$

An instrument with a weak correlation with the regressor gives an imprecise estimator. With a perfect correlation, then we get the precision of the LS estimator (which is precise, but perhaps not consistent).

### 5.4.2 2SLS

Suppose now that we have more instruments, $z_{t}$, than regressors, $x_{t}$. The IV method does not work since, there are then more equations than unknowns in (5.20). Instead, we can use the $2 S L S$ estimator. It has two steps. First, regress all elements in $x_{t}$ on all elements in $z_{t}$ with LS. Second, use the fitted values of $x_{t}$, denoted $\hat{x}_{t}$, as instruments in the IV method (use $\hat{x}_{t}$ in place of $z_{t}$ in the equations above). In can be shown that this is the most
efficient use of the information in $z_{t}$. The IV is clearly a special case of 2SLS (when $z_{t}$ has the same number of elements as $x_{t}$ ).

It is immediate from (5.22) that 2SLS is consistent under the same condiditons as IV since $\hat{x}_{t}$ is a linear function of the instruments, so $\operatorname{plim} \sum_{t=1}^{T} \hat{x}_{t} u_{t} / T=0$, if all the instruments are uncorrelated with $u_{t}$.

The name, 2SLS, comes from the fact that we get exactly the same result if we replace the second step with the following: regress $y_{t}$ on $\hat{x}_{t}$ with LS.

Example 9 (Supply equation with 2SLS.). With only one instrument, $A_{t}$, this is the same as Example 6, but presented in another way. First, regress $p_{t}$ on $A_{t}$

$$
p_{t}=\delta A_{t}+u_{t} \Rightarrow \operatorname{plim} \hat{\delta}_{L S}=\frac{\operatorname{Cov}\left(p_{t}, A_{t}\right)}{\operatorname{Var}\left(A_{t}\right)}=-\frac{1}{\beta-\gamma} \alpha
$$

Construct the predicted values as

$$
\hat{p}_{t}=\hat{\delta}_{L S} A_{t} .
$$

Second, regress $q_{t}$ on $\hat{p}_{t}$

$$
q_{t}=\gamma \hat{p}_{t}+e_{t}, \text { with } \operatorname{plim} \hat{\gamma}_{2 S L S}=\operatorname{plim} \frac{\widehat{\operatorname{Cov}}\left(q_{t}, \hat{p}_{t}\right)}{\widehat{\operatorname{Var}}\left(\hat{p}_{t}\right)} .
$$

Use $\hat{p}_{t}=\hat{\delta}_{L S} A_{t}$ and Slutsky's theorem

$$
\begin{aligned}
\operatorname{plim} \hat{\gamma}_{2 S L S} & =\frac{\operatorname{plim} \widehat{\operatorname{Cov}}\left(q_{t}, \hat{\delta}_{L S} A_{t}\right)}{\operatorname{plim} \widehat{\operatorname{Var}}\left(\hat{\delta}_{L S} A_{t}\right)} \\
& =\frac{\operatorname{Cov}\left(q_{t}, A_{t}\right) \operatorname{plim} \hat{\delta}_{L S}}{\operatorname{Var}\left(A_{t}\right) \operatorname{plim} \hat{\delta}_{L S}^{2}} \\
& =\frac{\left[-\frac{\gamma}{\beta-\gamma} \alpha \operatorname{Var}\left(A_{t}\right)\right]\left[-\frac{1}{\beta-\gamma} \alpha\right]}{\operatorname{Var}\left(A_{t}\right)\left[-\frac{1}{\beta-\gamma} \alpha\right]^{2}} \\
& =\gamma .
\end{aligned}
$$

Note that the trick here is to suppress some the movements in $p_{t}$. Only those movements that depend on $A_{t}$ (the observable shifts of the demand curve) are used. Movements in $p_{t}$ which are due to the unobservable demand and supply shocks are disregarded in $\hat{p}_{t}$. We
know from Example 2 that it is the supply shocks that make the LS estimate of the supply curve inconsistent. The IV method suppresses both them and the unobservable demand shock.

### 5.5 Hausman's Specification Test*

## Reference: Greene (2000) 9.5

This test is constructed to test if an efficient estimator (like LS) gives (approximately) the same estimate as a consistent estimator (like IV). If not, the efficient estimator is most likely inconsistent. It is therefore a way to test for the presence of endogeneity and/or measurement errors.

Let $\hat{\beta}_{e}$ be an estimator that is consistent and asymptotically efficient when the null hypothesis, $H_{0}$, is true, but inconsistent when $H_{0}$ is false. Let $\hat{\beta}_{c}$ be an estimator that is consistent under both $H_{0}$ and the alternative hypothesis. When $H_{0}$ is true, the asymptotic distribution is such that

$$
\begin{equation*}
\operatorname{Cov}\left(\hat{\beta}_{e}, \hat{\beta}_{c}\right)=\operatorname{Var}\left(\hat{\beta}_{e}\right) . \tag{5.27}
\end{equation*}
$$

Proof. Consider the estimator $\lambda \hat{\beta}_{c}+(1-\lambda) \hat{\beta}_{e}$, which is clearly consistent under $H_{0}$ since both $\hat{\beta}_{c}$ and $\hat{\beta}_{e}$ are. The asymptotic variance of this estimator is

$$
\lambda^{2} \operatorname{Var}\left(\hat{\beta}_{c}\right)+(1-\lambda)^{2} \operatorname{Var}\left(\hat{\beta}_{e}\right)+2 \lambda(1-\lambda) \operatorname{Cov}\left(\hat{\beta}_{c}, \hat{\beta}_{e}\right)
$$

which is minimized at $\lambda=0$ (since $\hat{\beta}_{e}$ is asymptotically efficient). The first order condition with respect to $\lambda$

$$
2 \lambda \operatorname{Var}\left(\hat{\beta}_{c}\right)-2(1-\lambda) \operatorname{Var}\left(\hat{\beta}_{e}\right)+2(1-2 \lambda) \operatorname{Cov}\left(\hat{\beta}_{c}, \hat{\beta}_{e}\right)=0
$$

should therefore be zero at $\lambda=0$ so

$$
\operatorname{Var}\left(\hat{\beta}_{e}\right)=\operatorname{Cov}\left(\hat{\beta}_{c}, \hat{\beta}_{e}\right) .
$$

(See Davidson (2000) 8.1)

This means that we can write

$$
\begin{align*}
\operatorname{Var}\left(\hat{\beta}_{e}-\hat{\beta}_{c}\right) & =\operatorname{Var}\left(\hat{\beta}_{e}\right)+\operatorname{Var}\left(\hat{\beta}_{c}\right)-2 \operatorname{Cov}\left(\hat{\beta}_{e}, \hat{\beta}_{c}\right) \\
& =\operatorname{Var}\left(\hat{\beta}_{c}\right)-\operatorname{Var}\left(\hat{\beta}_{e}\right) \tag{5.28}
\end{align*}
$$

We can use this to test, for instance, if the estimates from least squares ( $\hat{\beta}_{e}$, since LS is efficient if errors are iid normally distributed) and instrumental variable method ( $\hat{\beta}_{c}$, since consistent even if the true residuals are correlated with the regressors) are the same. In this case, $H_{0}$ is that the true residuals are uncorrelated with the regressors.

All we need for this test are the point estimates and consistent estimates of the variance matrices. Testing one of the coefficient can be done by a $t$ test, and testing all the parameters by a $\chi^{2}$ test

$$
\begin{equation*}
\left(\hat{\beta}_{e}-\hat{\beta}_{c}\right)^{\prime} \operatorname{Var}\left(\hat{\beta}_{e}-\hat{\beta}_{c}\right)^{-1}\left(\hat{\beta}_{e}-\hat{\beta}_{c}\right) \sim \chi^{2}(j) \tag{5.29}
\end{equation*}
$$

where $j$ equals the number of regressors that are potentially endogenous or measured with error. Note that the covariance matrix in (5.28) and (5.29) is likely to have a reduced rank, so the inverse needs to be calculated as a generalized inverse.

### 5.6 Tests of Overidentifying Restrictions in 2SLS*

When we use 2SLS, then we can test if instruments affect the dependent variable only via their correlation with the regressors. If not, something is wrong with the model since some relevant variables are excluded from the regression.

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## 6 Simulating the Finite Sample Properties

Reference: Greene (2000) 5.3
Additional references: Cochrane (2001) 15.2; Davidson and MacKinnon (1993) 21; Davison and Hinkley (1997); Efron and Tibshirani (1993) (bootstrapping, chap 9 in particular); and Berkowitz and Kilian (2000) (bootstrapping in time series models)

We know the small sample properties of regression coefficients in linear models with fixed regressors ( $X$ is non-stochastic) and iid normal error terms. Monte Carlo Simulations and bootstrapping are two common techniques used to understand the small sample properties when these conditions are not satisfied.

### 6.1 Monte Carlo Simulations in the Simplest Case

Monte Carlo simulations is essentially a way to generate many artificial (small) samples from a parameterized model and then estimating the statistics on each of those samples. The distribution of the statistics is then used as the small sample distribution of the estimator.

The following is an example of how Monte Carlo simulations could be done in the special case of a linear model for a scalar dependent variable

$$
\begin{equation*}
y_{t}=x_{t}^{\prime} \beta+u_{t} \tag{6.1}
\end{equation*}
$$

where $u_{t}$ is iid $N\left(0, \sigma^{2}\right)$ and $x_{t}$ is stochastic but independent of $u_{t \pm s}$ for all $s$. This means that $x_{t}$ cannot include lags of $y_{t}$.

Suppose we want to find the small sample distribution of a function of the estimate, $g(\hat{\beta})$. To do a Monte Carlo experiment, we need information on (i) $\beta$; (ii) the variance of $u_{t}, \sigma^{2}$; (iii) and a process for $x_{t}$.

The process for $x_{t}$ is typically estimated from the data on $x_{t}$. For instance, we could estimate the VAR system $x_{t}=A_{1} x_{t-1}+A_{2} x_{t-2}+e_{t}$. An alternative is to take an actual sample of $x_{t}$ and repeat it.

The values of $\beta$ and $\sigma^{2}$ are often a mix of estimation results and theory. In some
case, we simply take the point estimates. In other cases, we adjust the point estimates so that $g(\beta)=0$ holds, that is, so you simulate the model under the null hypothesis in order to study the size of asymptotic tests and to find valid critical values for small samples. Alternatively, you may simulate the model under an alternative hypothesis in order to study the power of the test using either critical values from either the asymptotic distribution or from a (perhaps simulated) small sample distribution.

To make it a bit concrete, suppose you want to use these simulations to get a $5 \%$ critical value for testing the null hypothesis $g(\beta)=0$. The Monte Carlo experiment follows these steps.

1. (a) Construct an artificial sample of the regressors (see above), $\left\{\tilde{x}_{t}\right\}_{t=1}^{T}$.
(b) Draw random numbers $\left\{\tilde{u}_{t}\right\}_{t=1}^{T}$ and use those together with the artificial sample of $\tilde{x}_{t}$ to calculate an artificial sample $\left\{\tilde{y}_{t}\right\}_{t=1}^{T}$ by using (6.1). Calculate an estimate $\hat{\beta}$ and record it along with the value of $g(\hat{\beta})$ and perhaps also the test statistics of the hypothesis that $g(\beta)=0$.
2. Repeat the previous steps $N(3000$, say $)$ times. The more times you repeat, the better is the approximation of the small sample distribution.
3. Sort your simulated $\hat{\beta}, g(\hat{\beta})$, and the test statistics in ascending order. For a onesided test (for instance, a chi-square test), take the $(0.95 N)$ th observations in these sorted vector as your 5\% critical values. For a two-sided test (for instance, a t-test), take the $(0.025 N)$ th and $(0.975 N)$ th observations as the $5 \%$ critical values. You can also record how many times the $5 \%$ critical values from the asymptotic distribution would reject a true null hypothesis.
4. You may also want to plot a histogram of $\hat{\beta}, g(\hat{\beta})$, and the test statistics to see if there is a small sample bias, and how the distribution looks like. Is it close to normal? How wide is it?
5. See Figure 6.1 for an example.

Remark 1 (Generating $N(\mu, \Sigma)$ random numbers) Suppose you want to draw an $n \times 1$ vector $\varepsilon_{t}$ of $N(\mu, \Sigma)$ variables. Use the Cholesky decomposition to calculate the lower triangular $P$ such that $\Sigma=P P^{\prime}$ (note that Gauss and MatLab returns $P^{\prime}$ instead of



Figure 6.1: Results from a Monte Carlo experiment of LS estimation of the AR coefficient. Data generated by an AR(1) process, 5000 simulations.
$P$ ). Draw $u_{t}$ from an $N(0, I)$ distribution (randn in MatLab, rndn in Gauss), and define $\varepsilon_{t}=\mu+P u_{t}$. Note that $\operatorname{Cov}\left(\varepsilon_{t}\right)=\mathrm{E} P u_{t} u_{t}^{\prime} P^{\prime}=P I P^{\prime}=\Sigma$.

### 6.2 Monte Carlo Simulations in More Complicated Cases*

### 6.2.1 When $x_{t}$ Includes Lags of $y_{t}$

If $x_{t}$ contains lags of $y_{t}$, then we must set up the simulations so that feature is preserved in every artificial sample that we create. For instance, suppose $x_{t}$ includes $y_{t-1}$ and another vector $z_{t}$ of variables which are independent of $u_{t \pm s}$ for all $s$. We can then generate an artificial sample as follows. First, create a sample $\left\{\tilde{z}_{t}\right\}_{t=1}^{T}$ by some time series model or by taking the observed sample itself (as we did with $x_{t}$ in the simplest case). Second, observation $t$ of $\left\{\tilde{x}_{t}, \tilde{y}_{t}\right\}$ is generated as

$$
\tilde{x}_{t}=\left[\begin{array}{c}
\tilde{y}_{t-1}  \tag{6.2}\\
\tilde{z}_{t}
\end{array}\right] \text { and } \tilde{y}_{t}=\tilde{x}_{t}^{\prime} \beta+u_{t}
$$

which is repeated for $t=1, \ldots, T$. We clearly need the initial value $y_{0}$ to start up the artificial sample, so one observation from the original sample is lost.

del. $R_{t}=0.9 f_{t}+\varepsilon_{t}$, where $\varepsilon_{\mathrm{t}}$ has a $t_{3}$ distribution
Kurtosis for $T=10100$ 1000: 46.9 6.14 .1
Rejection rates of abs(t-stat)>1.645: 0.16 0.10 0.10
Rejection rates of abs(t-stat)>1.96:0.10 0.05 0.06

Figure 6.2: Results from a Monte Carlo experiment with thick-tailed errors. The regressor is iid normally distributed. The errors have a $t_{3}$-distribution, 5000 simulations.

### 6.2.2 More Complicated Errors

It is straightforward to sample the errors from other distributions than the normal, for instance, a uniform distribution. Equipped with uniformly distributed random numbers, you can always (numerically) invert the cumulative distribution function (cdf) of any distribution to generate random variables from any distribution by using the probability transformation method. See Figure 6.2 for an example.

Remark 2 Let $X \sim U(0,1)$ and consider the transformation $Y=F^{-1}(X)$, where $F^{-1}()$ is the inverse of a strictly increasing cdf $F$, then $Y$ has the CDF F (). (Proof: follows from the lemma on change of variable in a density function.)

Example 3 The exponential cdf is $x=1-\exp (-\theta y)$ with inverse $y=-\ln (1-x) / \theta$ Draw $x$ from $U(0.1)$ and transform to $y$ to get an exponentially distributed variable.

It is more difficult to handle non-iid errors, for instance, heteroskedasticity and autocorrelation. We then need to model the error process and generate the errors from that model. For instance, if the errors are assumed to follow an $\operatorname{AR}(2)$ process, then we could estimate that process from the errors in (6.1) and then generate artificial samples of errors.

### 6.3 Bootstrapping in the Simplest Case

Bootstrapping is another way to do simulations, where we construct artificial samples by sampling from the actual data. The advantage of the bootstrap is then that we do not have to try to estimate the process of the errors and regressors as we must do in a Monte Carlo experiment. The real benefit of this is that we do not have to make any strong assumption about the distribution of the errors.

The bootstrap approach works particularly well when the errors are iid and independent of $x_{t-s}$ for all $s$. This means that $x_{t}$ cannot include lags of $y_{t}$. We here consider bootstrapping the linear model (6.1), for which we have point estimates (perhaps from LS) and fitted residuals. The procedure is similar to the Monte Carlo approach, except that the artificial sample is generated differently. In particular, Step 1 in the Monte Carlo simulation is replaced by the following:

1. Construct an artificial sample $\left\{\tilde{y}_{t}\right\}_{t=1}^{T}$ by

$$
\begin{equation*}
\tilde{y}_{t}=x_{t}^{\prime} \beta+\tilde{u}_{t}, \tag{6.3}
\end{equation*}
$$

where $\tilde{u}_{t}$ is drawn (with replacement) from the fitted residual and where $\beta$ is the point estimate. Calculate an estimate $\hat{\beta}$ and record it along with the value of $g(\hat{\beta})$ and perhaps also the test statistics of the hypothesis that $g(\beta)=0$.

### 6.4 Bootstrapping in More Complicated Cases*

### 6.4.1 Case 2: Errors are iid but Correlated With $x_{t+s}$

When $x_{t}$ contains lagged values of $y_{t}$, then we have to modify the approach in (6.3) since $\tilde{u}_{t}$ can become correlated with $x_{t}$. For instance, if $x_{t}$ includes $y_{t-1}$ and we happen to sample $\tilde{u}_{t}=\hat{u}_{t-1}$, then we get a non-zero correlation. The easiest way to handle this is as
in the Monte Carlo simulations: replace any $y_{t-1}$ in $x_{t}$ by $\tilde{y}_{t-1}$, that is, the corresponding observation in the artificial sample.

### 6.4.2 Case 3: Errors are Heteroskedastic but Uncorrelated with of $x_{t \pm s}$

Case 1 and 2 both draw errors randomly-based on the assumption that the errors are iid. Suppose instead that the errors are heteroskedastic, but still serially uncorrelated. We know that if the heteroskedastcity is related to the regressors, then the traditional LS covariance matrix is not correct (this is the case that White's test for heteroskedasticity tries to identify). It would then be wrong it pair $x_{t}$ with just any $\hat{u}_{s}$ since that destroys the relation between $x_{t}$ and the variance of $u_{t}$.

An alternative way of bootstrapping can then be used: generate the artificial sample by drawing (with replacement) pairs $\left(y_{s}, x_{s}\right)$, that is, we let the artificial pair in $t$ be $\left(\tilde{y}_{t}, \tilde{x}_{t}\right)=\left(x_{s}^{\prime} \hat{\beta}_{0}+\hat{u}_{s}, x_{s}\right)$ for some random draw of $s$ so we are always pairing the residual, $\hat{u}_{s}$, with the contemporaneous regressors, $x_{s}$. Note that is we are always sampling with replacement-otherwise the approach of drawing pairs would be just re-create the original data set. For instance, if the data set contains 3 observations, then artificial sample could be

$$
\left[\begin{array}{c}
\left(\tilde{y}_{1}, \tilde{x}_{1}\right) \\
\left(\tilde{y}_{2}, \tilde{x}_{2}\right) \\
\left(\tilde{y}_{3}, \tilde{x}_{3}\right)
\end{array}\right]=\left[\begin{array}{l}
\left(x_{2}^{\prime} \hat{\beta}_{0}+\hat{u}_{2}, x_{2}\right) \\
\left(x_{3}^{\prime} \hat{\beta}_{0}+\hat{u}_{3}, x_{3}\right) \\
\left(x_{3}^{\prime} \hat{\beta}_{0}+\hat{u}_{3}, x_{3}\right)
\end{array}\right]
$$

In contrast, when we sample (with replacement) $\hat{u}_{s}$, as we did above, then an artificial sample could be

$$
\left.\begin{array}{l}
\left(\tilde{y}_{1}, \tilde{x}_{1}\right) \\
\left(\tilde{y}_{2}, \tilde{x}_{2}\right) \\
\left(\tilde{y}_{3}, \tilde{x}_{3}\right)
\end{array}\right]=\left[\begin{array}{l}
\left(x_{1}^{\prime} \hat{\beta}_{0}+\hat{u}_{2}, x_{1}\right) \\
\left(x_{2}^{\prime} \hat{\beta}_{0}+\hat{u}_{1}, x_{2}\right) \\
\left(x_{3}^{\prime} \hat{\beta}_{0}+\hat{u}_{2}, x_{3}\right)
\end{array}\right]
$$

Davidson and MacKinnon (1993) argue that bootstrapping the pairs ( $y_{s}, x_{s}$ ) makes little sense when $x_{s}$ contains lags of $y_{s}$, since there is no way to construct lags of $y_{s}$ in the bootstrap. However, what is important for the estimation is sample averages of various functions of the dependent and independent variable within a period—not how the line up over time (provided the assumption of no autocorrelation of the residuals is true)

### 6.4.3 Other Approaches

There are many other ways to do bootstrapping. For instance, we could sample the regressors and residuals independently of each other and construct an artificial sample of the dependent variable $\tilde{y}_{t}=\tilde{x}_{t}^{\prime} \hat{\beta}+\tilde{u}_{t}$. This clearly makes sense if the residuals and regressors are independent of each other and errors are iid. In that case, the advantage of this approach is that we do not keep the regressors fixed

### 6.4.4 Serially Dependent Errors

It is quite hard to handle the case when the errors are serially dependent, since we must the sample in such a way that we do not destroy the autocorrelation structure of the data. A common approach is to fit a model for the residuals, for instance, an $\operatorname{AR}(1)$, and then bootstrap the (hopefully iid) innovations to that process.
Another approach amounts to resampling of blocks of data. For instance, suppose the sample has 10 observations, and we decide to create blocks of 3 observations. The first block is ( $\hat{u}_{1}, \hat{u}_{2}, \hat{u}_{3}$ ), the second block is ( $\hat{u}_{2}, \hat{u}_{3}, \hat{u}_{4}$ ), and so forth until the last block, ( $\hat{u}_{8}, \hat{u}_{9}, \hat{u}_{10}$ ). If we need a sample of length $3 \tau$, say, then we simply draw $\tau$ of those block randomly (with replacement) and stack them to form a longer series. To handle end point effects (so that all data points have the same probability to be drawn), we also create blocks by "wrapping" the data around a circle. In practice, this means that we add a the following blocks: $\left(\hat{u}_{10}, \hat{u}_{1}, \hat{u}_{2}\right)$ and $\left(\hat{u}_{9}, \hat{u}_{10}, \hat{u}_{1}\right)$. An alternative approach is to have non-overlapping blocks. See Berkowitz and Kilian (2000) for some other recent methods.

## Bibliography

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## 7 GMM

References: Greene (2000) 4.7 and 11.5-6
Additional references: Hayashi (2000) 3-4; Verbeek (2000) 5; Hamilton (1994) 14; Ogaki (1993), Johnston and DiNardo (1997) 10; Harris and Matyas (1999); Pindyck and Rubinfeld (1997) Appendix 10.1; Cochrane (2001) 10-11

### 7.1 Method of Moments

Let $m\left(x_{t}\right)$ be a $k \times 1$ vector valued continuous function of a stationary process, and let the probability limit of the mean of $m$ (.) be a function $\gamma$ (.) of a $k \times 1$ vector $\beta$ of parameters. We want to estimate $\beta$. The method of moments (MM, not yet generalized to GMM) estimator is obtained by replacing the probability limit with the sample mean and solving the system of $k$ equations

$$
\begin{equation*}
\frac{1}{T} \sum_{t=1}^{T} m\left(x_{t}\right)-\gamma(\beta)=\mathbf{0}_{k \times 1} \tag{7.1}
\end{equation*}
$$

## for the parameters $\beta$.

It is clear that this is a consistent estimator of $\beta$ if $\gamma$ is continuous. (Proof: the sample mean is a consistent estimator of $\gamma($.$) , and by Slutsky's theorem \operatorname{plim} \gamma(\hat{\beta})=\gamma(\operatorname{plim} \hat{\beta})$ if $\gamma$ is a continuous function.)

Example 1 (MM for the variance of a variable.) The MM condition $\frac{1}{T} \sum_{t=1}^{T} x_{t}^{2}-\sigma^{2}=0$ gives the usual MLE of the sample variance, assuming $E x_{t}=0$.

Example 2 (MM for an MA(1).) For an $M A(1)$, $y_{t}=\epsilon_{t}+\theta \epsilon_{t-1}$, we have

$$
\begin{array}{clcc}
E y_{t}^{2} & = & E\left(\epsilon_{t}+\theta \epsilon_{t-1}\right)^{2} & =\sigma_{\epsilon}^{2}\left(1+\theta^{2}\right) \\
E\left(y_{t} y_{t-1}\right) & =E\left[\left(\epsilon_{t}+\theta \epsilon_{t-1}\right)\left(\epsilon_{t-1}+\theta \epsilon_{t-2}\right)\right] & = & \sigma_{\epsilon}^{2} \theta
\end{array}
$$

The moment conditions could therefore be

$$
\left[\begin{array}{c}
\frac{1}{T} \sum_{t=1}^{T} y_{t}^{2}-\sigma_{\epsilon}^{2}\left(1+\theta^{2}\right) \\
\frac{1}{T} \sum_{t=1}^{T} y_{t} y_{t-1}-\sigma_{\epsilon}^{2} \theta
\end{array}\right]=\left[\begin{array}{l}
0 \\
0
\end{array}\right],
$$

which allows us to estimate $\theta$ and $\sigma^{2}$.

### 7.2 Generalized Method of Moments

GMM extends MM by allowing for more orthogonality conditions than parameters. This could, for instance, increase efficiency and/or provide new aspects which can be tested.

Many (most) traditional estimation methods, like LS, IV, and MLE are special cases of GMM. This means that the properties of GMM are very general, and therefore fairly difficult to prove.

### 7.3 Moment Conditions in GMM

Suppose we have $q$ (unconditional) moment conditions,

$$
\begin{align*}
\mathrm{E} m\left(w_{t}, \beta_{0}\right) & =\left[\begin{array}{c}
\mathrm{E} m_{1}\left(w_{t}, \beta_{0}\right) \\
\vdots \\
\mathrm{E} m_{q}\left(w_{t}, \beta_{0}\right)
\end{array}\right] \\
& =\mathbf{0}_{q \times 1} \tag{7.2}
\end{align*}
$$

from which we want to estimate the $k \times 1(k \leq q)$ vector of parameters, $\beta$. The true values are $\beta_{0}$. We assume that $w_{t}$ is a stationary and ergodic (vector) process (otherwise the sample means does not converge to anything meaningful as the sample size increases). The sample averages, or "sample moment conditions," evaluated at some value of $\beta$, are

$$
\begin{equation*}
\bar{m}(\beta)=\frac{1}{T} \sum_{t=1}^{T} m\left(w_{t}, \beta\right) \tag{7.3}
\end{equation*}
$$

The sample average $\bar{m}(\beta)$ is a vector of functions of random variables, so they are random variables themselves and depend on the sample used. It will later be interesting to calculate the variance of $\bar{m}(\beta)$. Note that $\bar{m}\left(\beta_{1}\right)$ and $\bar{m}\left(\beta_{2}\right)$ are sample means obtained
by using two different parameter vectors, but on the same sample of data.
Example 3 (Moments conditions for IV/2SLS.) Consider the linear model $y_{t}=x_{t}^{\prime} \beta_{0}+u_{t}$, where $x_{t}$ and $\beta$ are $k \times 1$ vectors. Let $z_{t}$ be a $q \times 1$ vector, with $q \geq k$. The moment conditions and their sample analogues are

$$
\mathbf{0}_{q \times 1}=E z_{t} u_{t}=E\left[z_{t}\left(y_{t}-x_{t}^{\prime} \beta_{0}\right)\right], \text { and } \bar{m}(\beta)=\frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \beta\right),
$$

(or $Z^{\prime}(Y-X \beta) / T$ in matrix form). Let $q=k$ to get IV; let $z_{t}=x_{t}$ to get $L S$.
Example 4 (Moments conditions for MLE.) The maximum likelihood estimator maximizes the $\log$ likelihood function, $\frac{1}{T} \Sigma_{t=1}^{T} \ln L\left(w_{t} ; \beta\right)$, which requires $\frac{1}{T} \Sigma_{t=1}^{T} \partial \ln L\left(w_{t} ; \beta\right) / \partial \beta=$ 0. A key regularity condition for the $M L E$ is that $E \partial \ln L\left(w_{t} ; \beta_{0}\right) / \partial \beta=0$, which is just like a GMM moment condition.

### 7.3.1 Digression: From Conditional to Unconditional Moment Conditions

Suppose we are instead given conditional moment restrictions

$$
\begin{equation*}
\mathrm{E}\left[u\left(x_{t}, \beta_{0}\right) \mid z_{t}\right]=\mathbf{0}_{m \times 1}, \tag{7.4}
\end{equation*}
$$

where $z_{t}$ is a vector of conditioning (predetermined) variables. We want to transform this to unconditional moment conditions.

Remark $5(E(u \mid z)=0$ versus Euz $=0$.) For any random variables $u$ and $z$,

$$
\operatorname{Cov}(z, u)=\operatorname{Cov}[z, E(u \mid z)]
$$

The condition $E(u \mid z)=0$ then implies $\operatorname{Cov}(z, u)=0$. Recall that $\operatorname{Cov}(z, u)=E z u-E z E u$, and that $E(u \mid z)=0$ implies that $E u=0$ (by iterated expectations). We therefore get that

$$
E(u \mid z)=0 \Rightarrow\left[\begin{array}{c}
\operatorname{Cov}(z, u)=0 \\
E u=0
\end{array}\right] \Rightarrow E u z=0 .
$$

Example 6 (Euler equation for optimal consumption.) The standard Euler equation for optimal consumption choice which with isoelastic utility $U\left(C_{t}\right)=C_{t}^{1-\gamma} /(1-\gamma)$ is

$$
E\left[\left.R_{t+1} \beta\left(\frac{C_{t+1}}{C_{t}}\right)^{-\gamma}-1 \right\rvert\, \Omega_{t}\right]=0
$$

### 7.4 The Optimization Problem in GMM

### 7.4.1 The Loss Function

The GMM estimator $\hat{\beta}$ minimizes the weighted quadratic form

$$
J=\left[\begin{array}{l}
\bar{m}_{1}(\beta) \\
\vdots \\
\vdots \\
\bar{m}_{q}(\beta)
\end{array}\right]^{\prime}\left[\begin{array}{cccc}
W_{11} & \cdots & \cdots & W_{1 q} \\
\vdots & \ddots & & \vdots \\
\vdots & & \ddots & \vdots \\
W_{1 q} & \cdots & \cdots & W_{q q}
\end{array}\right]\left[\begin{array}{l}
\bar{m}_{1}(\beta) \\
\vdots \\
\vdots \\
\bar{m}_{q}(\beta)
\end{array}\right]
$$

$$
=\bar{m}(\beta)^{\prime} W \bar{m}(\beta),
$$

where $\bar{m}(\beta)$ is the sample average of $m\left(w_{t}, \beta\right)$ given by (7.3), and where $W$ is some $q \times q$ symmetric positive definite weighting matrix. (We will soon discuss a good choice of weighting matrix.) There are $k$ parameters in $\beta$ to estimate, and we have $q$ moment conditions in $\bar{m}(\beta)$. We therefore have $q-k$ overidentifying moment restrictions.

With $q=k$ the model is exactly identified (as many equations as unknown), and it should be possible to set all $q$ sample moment conditions to zero by a choosing the $k=q$ parameters. It is clear that the choice of the weighting matrix has no effect in this case since $\bar{m}(\hat{\beta})=0$ at the point estimates $\hat{\beta}$.

Example 8 (Simple linear regression.) Consider the model

$$
\begin{equation*}
y_{t}=x_{t} \beta_{0}+u_{t} \tag{7.9}
\end{equation*}
$$

where $y_{t}$ and $x_{t}$ are zero mean scalars. The moment condition and loss function are

$$
\begin{aligned}
\bar{m}(\beta) & =\frac{1}{T} \sum_{t=1}^{T} x_{t}\left(y_{t}-x_{t} \beta\right) \text { and } \\
J & =W\left[\frac{1}{T} \sum_{t=1}^{T} x_{t}\left(y_{t}-x_{t} \beta\right)\right]^{2},
\end{aligned}
$$

so the scalar $W$ is clearly irrelevant in this case.

Example 9 (IV/2SLS method continued.) From Example 3, we note that the loss function
for the IV/2SLS method is

$$
\bar{m}(\beta)^{\prime} W \bar{m}(\beta)=\left[\frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \beta\right)\right]^{\prime} W\left[\frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \beta\right)\right] .
$$

When $q=k$, then the model is exactly identified, so the estimator could actually be found by setting all moment conditions to zero. We then get the IV estimator

$$
\begin{aligned}
\mathbf{0} & =\frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \hat{\beta}_{I V}\right) \text { or } \\
\hat{\beta}_{I V} & =\left(\frac{1}{T} \sum_{t=1}^{T} z_{t} x_{t}^{\prime}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} z_{t} y_{t} \\
& =\hat{\Sigma}_{z x}^{-1} \hat{\Sigma}_{z y}
\end{aligned}
$$

where $\hat{\Sigma}_{z x}=\Sigma_{t=1}^{T} z_{t} x_{t}^{\prime} / T$ and similarly for the other second moment matrices. Let $z_{t}=$ $x_{t}$ to get $L S$

$$
\hat{\beta}_{L S}=\hat{\Sigma}_{x x}^{-1} \hat{\Sigma}_{x y} .
$$

### 7.4.2 First Order Conditions

Remark 10 (Matrix differentiation of non-linear functions.) Let the vector $y_{n \times 1}$ be a function of the vector $x_{m \times 1}$

$$
\left[\begin{array}{l}
y_{1} \\
\vdots \\
y_{n}
\end{array}\right]=f(x)=\left[\begin{array}{l}
f_{1}(x) \\
\vdots \\
f_{n}(x)
\end{array}\right]
$$

Then, $\partial y / \partial x^{\prime}$ is an $n \times m$ matrix

$$
\frac{\partial y}{\partial x^{\prime}}=\left[\begin{array}{l}
\frac{\partial f_{1}(x)}{\partial x^{\prime}} \\
\vdots \\
\frac{\partial f_{1}(x)}{\partial x^{\prime}}
\end{array}\right]=\left[\begin{array}{lll}
\frac{\partial f_{1}(x)}{\partial x_{1}} & \cdots & \frac{\partial f_{1}(x)}{\partial x_{m}} \\
\vdots & & \vdots \\
\frac{\partial f_{n}(x)}{\partial x_{1}} & \cdots & \frac{\partial f_{n}(x)}{\partial x_{m}}
\end{array}\right]
$$

(Note that the notation implies that the derivatives of the first element in $y$, denoted $y_{1}$, with respect to each of the elements in $x^{\prime}$ are found in the first row of $\partial y / \partial x^{\prime}$. A rule to help memorizing the format of $\partial y / \partial x^{\prime}: y$ is a column vector and $x^{\prime}$ is a row vector.)

Remark 11 When $y=A x$ where $A$ is an $n \times m$ matrix, then $f_{i}(x)$ in Remark 10 is a linear function. We then get $\partial y / \partial x^{\prime}=\partial(A x) / \partial x^{\prime}=A$.

Remark 12 As a special case of the previous remark $y=z^{\prime} x$ where both $z$ and $x$ are vectors. Then $\partial\left(z^{\prime} x\right) / \partial x^{\prime}=z^{\prime}\left(\right.$ since $z^{\prime}$ plays the role of $\left.A\right)$.

Remark 13 (Matrix differentiation of quadratic forms.) Let $x_{n \times 1}, f(x)_{m \times 1}$, and $A_{m \times m}$ symmetric. Then

$$
\frac{\partial f(x)^{\prime} A f(x)}{\partial x}=2\left(\frac{\partial f(x)}{\partial x^{\prime}}\right)^{\prime} A f(x)
$$

Remark 14 If $f(x)=x$, then $\partial f(x) / \partial x^{\prime}=I$, so $\partial\left(x^{\prime} A x\right) / \partial x=2 A x$.
The $k$ first order conditions for minimizing the GMM loss function in (7.8) with respect to the $k$ parameters are that the partial derivatives with respect to $\beta$ equal zero at the estimate, $\hat{\beta}$,


$$
\begin{equation*}
=\underbrace{\left(\frac{\partial \bar{m}(\hat{\beta})}{\partial \beta^{\prime}}\right)^{\prime}}_{k \times q} \underbrace{W}_{q \times q} \underbrace{\bar{m}(\hat{\beta})}_{q \times 1} . \tag{7.10}
\end{equation*}
$$

We can solve for the GMM estimator, $\hat{\beta}$, from (7.11). This set of equations must often be solved by numerical methods, except in linear models (the moment conditions are linear functions of the parameters) where we can find analytical solutions by matrix inversion.

Example 15 (First order conditions of simple linear regression.) The first order conditions of the loss function in Example 8 is

$$
\begin{aligned}
0 & =\frac{d}{d \beta} W\left[\frac{1}{T} \sum_{t=1}^{T} x_{t}\left(y_{t}-x_{t} \hat{\beta}\right)\right]^{2} \\
& =\left[-\frac{1}{T} \sum_{t=1}^{T} x_{t}^{2}\right] W\left[\frac{1}{T} \sum_{t=1}^{T} x_{t}\left(y_{t}-x_{t} \hat{\beta}\right)\right], \text { or } \\
\hat{\beta} & =\left(\frac{1}{T} \sum_{t=1}^{T} x_{t}^{2}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} x_{t} y_{t} .
\end{aligned}
$$

Example 16 (First order conditions of IV/2SLS.) The first order conditions corresponding to (7.11) of the loss function in Example 9 (when $q \geq k$ ) are

$$
\begin{aligned}
\mathbf{0}_{k \times 1} & =\left[\frac{\partial \bar{m}(\hat{\beta})}{\partial \beta^{\prime}}\right]^{\prime} W \bar{m}(\hat{\beta}) \\
& =\left[\frac{\partial}{\partial \beta^{\prime}} \frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \hat{\beta}\right)\right]^{\prime} W \frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \hat{\beta}\right) \\
& =\left[-\frac{1}{T} \sum_{t=1}^{T} z_{t} x_{t}^{\prime}\right]^{\prime} W \frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \hat{\beta}\right) \\
& =-\hat{\Sigma}_{x z} W\left(\hat{\Sigma}_{z y}-\hat{\Sigma}_{z x} \hat{\beta}\right) .
\end{aligned}
$$

We can solve for $\hat{\beta}$ from the first order conditions as

$$
\hat{\beta}_{2 S L S}=\left(\hat{\Sigma}_{x z} W \hat{\Sigma}_{z x}\right)^{-1} \hat{\Sigma}_{x z} W \hat{\Sigma}_{z y}
$$

When $q=k$, then the first order conditions can be premultiplied with $\left(\hat{\Sigma}_{x z} W\right)^{-1}$, since $\hat{\Sigma}_{x z} W$ is an invertible $k \times k$ matrix in this case, to give

$$
\mathbf{0}_{k \times 1}=\hat{\Sigma}_{z y}-\hat{\Sigma}_{z x} \hat{\beta}, \text { so } \hat{\beta}_{I V}=\hat{\Sigma}_{z x}^{-1} \hat{\Sigma}_{z y}
$$

This shows that the first order conditions are just the same as the sample moment conditions, which can be made to hold exactly since there are as many parameters as there are equations.

### 7.5 Asymptotic Properties of GMM

We know very little about the general small sample properties, including bias, of GMM We therefore have to rely either on simulations (Monte Carlo or bootstrap) or on the asymptotic results. This section is about the latter.

GMM estimates are typically consistent and normally distributed, even if the series $m\left(w_{t}, \beta\right)$ in the moment conditions (7.3) are serially correlated and heteroskedasticprovided $w_{t}$ is a stationary and ergodic process. The reason is essentially that the estimators are (at least as a first order approximation) linear combinations of sample means which typically are consistent (LLN) and normally distributed (CLT). More about that later. The proofs are hard, since the GMM is such a broad class of estimators. This section discusses, in an informal way, how we can arrive at those results.

### 7.5.1 Consistency

Sample moments are typically consistent, so $\operatorname{plim} m(\beta)=\mathrm{E} m\left(w_{t}, \beta\right)$. This must hold at any parameter vector in the relevant space (for instance, those inducing stationarity and variances which are strictly positive). Then, if the moment conditions (7.2) are true only at the true parameter vector, $\beta_{0}$, (otherwise the parameters are "unidentified") and that they are continuous in $\beta$, then GMM is consistent. The idea is thus that GMM asymptotically solves

$$
\begin{aligned}
\mathbf{0}_{q \times 1} & =\operatorname{plim} \bar{m}(\hat{\beta}) \\
& =\mathrm{E} m\left(w_{t}, \hat{\beta}\right),
\end{aligned}
$$

which only holds at $\hat{\beta}=\beta_{0}$. Note that this is an application of Slutsky's theorem.

Remark 17 (Slutsky's theorem.) If $\left\{x_{T}\right\}$ is a sequence of random matrices such that $\operatorname{plim} x_{T}=x$ and $g\left(x_{T}\right)$ a continuous function, then $\operatorname{plim} g\left(x_{T}\right)=g(x)$.

Example 18 (Consistency of $2 S L S$.) By using $y_{t}=x_{t}^{\prime} \beta_{0}+u_{t}$, the first order conditions
in Example 16 can be rewritten

$$
\begin{aligned}
\mathbf{0}_{k \times 1} & =\hat{\Sigma}_{x z} W \frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \hat{\beta}\right) \\
& =\hat{\Sigma}_{x z} W \frac{1}{T} \sum_{t=1}^{T} z_{t}\left[u_{t}+x_{t}^{\prime}\left(\beta_{0}-\hat{\beta}\right)\right] \\
& =\hat{\Sigma}_{x z} W \hat{\Sigma}_{z u}+\hat{\Sigma}_{x z} W \hat{\Sigma}_{z x}\left(\beta_{0}-\hat{\beta}\right) .
\end{aligned}
$$

Take the probability limit

$$
\mathbf{0}_{k \times 1}=\operatorname{plim} \hat{\Sigma}_{x z} W \operatorname{plim} \hat{\Sigma}_{z u}+\operatorname{plim} \hat{\Sigma}_{x z} W \operatorname{plim} \hat{\Sigma}_{z x}\left(\beta_{0}-\operatorname{plim} \hat{\beta}\right) .
$$

In most cases, $\operatorname{plim} \hat{\Sigma}_{x z}$ is some matrix of constants, and $\operatorname{plim} \hat{\Sigma}_{z u}=\mathrm{E} z_{t} u_{t}=\mathbf{0}_{q \times 1}$. It then follows that $\operatorname{plim} \hat{\beta}=\beta_{0}$. Note that the whole argument relies on that the moment condition, $E z_{t} u_{t}=\mathbf{0}_{q \times 1}$, is true. If it is not, then the estimator is inconsistent. For instance, when the instruments are invalid (correlated with the residuals) or when we use $L S\left(z_{t}=x_{t}\right)$ when there are measurement errors or in a system of simultaneous equations.

### 7.5.2 Asymptotic Normality

To give the asymptotic distribution of $\sqrt{T}\left(\hat{\beta}-\beta_{0}\right)$, we need to define three things. (As usual, we also need to scale with $\sqrt{T}$ to get a non-trivial asymptotic distribution; the asymptotic distribution of $\hat{\beta}-\beta_{0}$ is a spike at zero.) First, let $S_{0}$ (a $q \times q$ matrix) denote the asymptotic covariance matrix (as sample size goes to infinity) of $\sqrt{T}$ times the sample moment conditions evaluated at the true parameters

$$
\begin{align*}
S_{0} & =\mathrm{ACov}\left[\sqrt{T} \bar{m}\left(\beta_{0}\right)\right]  \tag{7.12}\\
& =\mathrm{ACov}\left[\frac{1}{\sqrt{T}} \sum_{t=1}^{T} m\left(w_{t}, \beta_{0}\right)\right], \tag{7.13}
\end{align*}
$$

where we use the definition of $\bar{m}\left(\beta_{0}\right)$ in (7.3). (To estimate $S_{0}$ it is important to recognize that it is a scaled sample average.) Second, let $D_{0}$ (a $q \times k$ matrix) denote the probability limit of the gradient of the sample moment conditions with respect to the parameters,
evaluated at the true parameters

$$
\begin{equation*}
D_{0}=\operatorname{plim} \frac{\partial \bar{m}\left(\beta_{0}\right)}{\partial \beta^{\prime}} \tag{7.14}
\end{equation*}
$$

Note that a similar gradient, but evaluated at $\hat{\beta}$, also shows up in the first order conditions (7.11). Third, let the weighting matrix be the inverse of the covariance matrix of the moment conditions (once again evaluated at the true parameters)

$$
\begin{equation*}
W=S_{0}^{-1} \tag{7.15}
\end{equation*}
$$

It can be shown that this choice of weighting matrix gives the asymptotically most efficient estimator for a given set of orthogonality conditions. For instance, in 2SLS, this means a given set of instruments and (7.15) then shows only how to use these instruments in the most efficient way. Of course, another set of instruments might be better (in the sense of giving a smaller $\operatorname{Cov}(\hat{\beta}))$.

With the definitions in (7.12) and (7.14) and the choice of weighting matrix in (7.15) and the added assumption that the rank of $D_{0}$ equals $k$ (number of parameters) then we can show (under fairly general conditions) that

$$
\begin{equation*}
\sqrt{T}\left(\hat{\beta}-\beta_{0}\right) \xrightarrow{d} N\left(\mathbf{0}_{k \times 1}, V\right), \text { where } V=\left(D_{0}^{\prime} S_{0}^{-1} D_{0}\right)^{-1} . \tag{7.16}
\end{equation*}
$$

This holds also when the model is exactly identified, so we really do not use any weighting matrix.

To prove this note the following.
Remark 19 (Continuous mapping theorem.) Let the sequences of random matrices $\left\{x_{T}\right\}$ and $\left\{y_{T}\right\}$, and the non-random matrix $\left\{a_{T}\right\}$ be such that $x_{T} \xrightarrow{d} x, y_{T} \xrightarrow{p} y$, and $a_{T} \rightarrow a$ (a traditional limit). Let $g\left(x_{T}, y_{T}, a_{T}\right)$ be a continuous function. Then $g\left(x_{T}, y_{T}, a_{T}\right) \xrightarrow{d}$ $g(x, y, a)$. Either of $y_{T}$ and $a_{T}$ could be irrelevant in g. (See Mittelhammer (1996) 5.3.)

Example 20 For instance, the sequences in Remark 19 could be $x_{T}=\sqrt{T} \Sigma_{t=}^{T} w_{t} / T$, the scaled sample average of a random variable $w_{t} ; y_{T}=\Sigma_{t=}^{T} w_{t}^{2} / T$, the sample second moment; and $a_{T}=\Sigma_{t=1}^{T} 0.7^{t}$.

Remark 21 From the previous remark: if $x_{T} \xrightarrow{d} x$ (a random variable) and $\operatorname{plim} Q_{T}=$ $Q$ (a constant matrix), then $Q_{T} x_{T} \xrightarrow{d} Q x$.

Proof. (The asymptotic distribution (7.16). Sketch of proof.) This proof is essentially an application of the delta rule. By the mean-value theorem the sample moment condition evaluated at the GMM estimate, $\hat{\beta}$, is

$$
\begin{equation*}
\bar{m}(\hat{\beta})=\bar{m}\left(\beta_{0}\right)+\frac{\partial \bar{m}\left(\beta_{1}\right)}{\partial \beta^{\prime}}\left(\hat{\beta}-\beta_{0}\right) \tag{7.17}
\end{equation*}
$$

for some values $\beta_{1}$ between $\hat{\beta}$ and $\beta_{0}$. (This point is different for different elements in $\bar{m}$.) Premultiply with $\left[\partial \bar{m}(\hat{\beta}) / \partial \beta^{\prime}\right]^{\prime} W$. By the first order condition (7.11), the left hand side is then zero, so we have

$$
\begin{equation*}
\mathbf{0}_{k \times 1}=\left(\frac{\partial \bar{m}(\hat{\beta})}{\partial \beta^{\prime}}\right)^{\prime} W \bar{m}\left(\beta_{0}\right)+\left(\frac{\partial \bar{m}(\hat{\beta})}{\partial \beta^{\prime}}\right)^{\prime} W \frac{\partial \bar{m}\left(\beta_{1}\right)}{\partial \beta^{\prime}}\left(\hat{\beta}-\beta_{0}\right) . \tag{7.18}
\end{equation*}
$$

Multiply with $\sqrt{T}$ and solve as

$$
\begin{equation*}
\sqrt{T}\left(\hat{\beta}-\beta_{0}\right)=\underbrace{-\left[\left(\frac{\partial \bar{m}(\hat{\beta})}{\partial \beta^{\prime}}\right)^{\prime} W \frac{\partial \bar{m}\left(\beta_{1}\right)}{\partial \beta^{\prime}}\right]^{-1}\left(\frac{\partial \bar{m}(\hat{\beta})}{\partial \beta^{\prime}}\right)^{\prime} W}_{\Gamma} \sqrt{T} \bar{m}\left(\beta_{0}\right) . \tag{7.19}
\end{equation*}
$$

If

$$
\operatorname{plim} \frac{\partial \bar{m}(\hat{\beta})}{\partial \beta^{\prime}}=\frac{\partial \bar{m}\left(\beta_{0}\right)}{\partial \beta^{\prime}}=D_{0}, \text { then } \operatorname{plim} \frac{\partial \bar{m}\left(\beta_{1}\right)}{\partial \beta^{\prime}}=D_{0},
$$

since $\beta_{1}$ is between $\beta_{0}$ and $\hat{\beta}$. Then

$$
\begin{equation*}
\operatorname{plim} \Gamma=-\left(D_{0}^{\prime} W D_{0}\right)^{-1} D_{0}^{\prime} W \tag{7.20}
\end{equation*}
$$

The last term in (7.19), $\sqrt{T} \bar{m}\left(\beta_{0}\right)$, is $\sqrt{T}$ times a vector of sample averages, so by a CLT it converges in distribution to $\mathrm{N}\left(0, S_{0}\right)$, where $S_{0}$ is defined as in (7.12). By the rules of limiting distributions (see Remark 19) we then have that

$$
\begin{aligned}
& \sqrt{T}\left(\hat{\beta}-\beta_{0}\right) \xrightarrow{d} \operatorname{plim} \Gamma \times \text { something that is } N\left(0, S_{0}\right), \text { that is, } \\
& \sqrt{T}\left(\hat{\beta}-\beta_{0}\right) \xrightarrow{d} N\left[\mathbf{0}_{k \times 1},(\operatorname{plim} \Gamma) S_{0}\left(\operatorname{plim} \Gamma^{\prime}\right)\right] .
\end{aligned}
$$

The covariance matrix is then

$$
\begin{align*}
\operatorname{ACov}\left[\sqrt{T}\left(\hat{\beta}-\beta_{0}\right)\right] & =(\operatorname{plim} \Gamma) S_{0}\left(\operatorname{plim} \Gamma^{\prime}\right) \\
& =\left(D_{0}^{\prime} W D_{0}\right)^{-1} D_{0}^{\prime} W S_{0}\left[\left(D_{0}^{\prime} W D_{0}\right)^{-1} D_{0}^{\prime} W\right]^{\prime}  \tag{7.21}\\
& =\left(D_{0}^{\prime} W D_{0}\right)^{-1} D_{0}^{\prime} W S_{0} W^{\prime} D_{0}\left(D_{0}^{\prime} W D_{0}\right)^{-1} \tag{7.22}
\end{align*}
$$

If $W=W^{\prime}=S_{0}^{-1}$, then this expression simplifies to (7.16). (See, for instance, Hamilton (1994) 14 (appendix) for more details.)

It is straightforward to show that the difference between the covariance matrix in (7.22) and $\left(D_{0}^{\prime} S_{0}^{-1} D_{0}\right)^{-1}$ (as in (7.16)) is a positive semi-definite matrix: any linear combination of the parameters has a smaller variance if $W=S_{0}^{-1}$ is used as the weighting matrix

All the expressions for the asymptotic distribution are supposed to be evaluated at the true parameter vector $\beta_{0}$, which is unknown. However, $D_{0}$ in (7.14) can be estimated by $\partial \bar{m}(\hat{\beta}) / \partial \beta^{\prime}$, where we use the point estimate instead of the true value of the parameter vector. In practice, this means plugging in the point estimates into the sample moment conditions and calculate the derivatives with respect to parameters (for instance, by a numerical method).

Similarly, $S_{0}$ in (7.13) can be estimated by, for instance, Newey-West's estimator of $\operatorname{Cov}[\sqrt{T} \bar{m}(\hat{\beta})]$, once again using the point estimates in the moment conditions.

### 7.6 Summary of GMM

$$
\begin{aligned}
& \text { Economic model : } \operatorname{E} m\left(w_{t}, \beta_{0}\right)=\mathbf{0}_{q \times 1}, \beta \text { is } k \times 1 \\
& \text { Sample moment conditions : } \bar{m}(\beta)=\frac{1}{T} \sum_{t=1}^{T} m\left(w_{t}, \beta\right) \\
& \text { Loss function : } J=\bar{m}(\beta)^{\prime} W \bar{m}(\beta) \\
& \text { First order conditions : } \mathbf{0}_{k \times 1}=\frac{\partial \bar{m}(\hat{\beta})^{\prime} W \bar{m}(\hat{\beta})}{\partial \beta}=\left(\frac{\partial \bar{m}(\hat{\beta})}{\partial \beta^{\prime}}\right)^{\prime} W \bar{m}(\hat{\beta}) \\
& \text { Consistency : } \hat{\beta} \text { is typically consistent if } \operatorname{Em}\left(w_{t}, \beta_{0}\right)=\mathbf{0} \\
& \text { Define : } S_{0}=\operatorname{Cov}\left[\sqrt{T} \bar{m}\left(\beta_{0}\right)\right] \text { and } D_{0}=\operatorname{plim} \frac{\partial \bar{m}\left(\beta_{0}\right)}{\partial \beta^{\prime}} \\
& \text { Choose: } W=S_{0}^{-1} \\
& \text { Asymptotic distribution : } \sqrt{T}\left(\hat{\beta}-\beta_{0}\right) \xrightarrow{d} N\left(\mathbf{0}_{k \times 1}, V\right) \text {, where } V=\left(D_{0}^{\prime} S_{0}^{-1} D_{0}\right)^{-1}
\end{aligned}
$$

### 7.7 Efficient GMM and Its Feasible Implementation

The efficient GMM (remember: for a given set of moment conditions) requires that we use $W=S_{0}^{-1}$, which is tricky since $S_{0}$ should be calculated by using the true (unknown) parameter vector. However, the following two-stage procedure usually works fine:

- First, estimate model with some (symmetric and positive definite) weighting matrix. The identity matrix is typically a good choice for models where the moment conditions are of the same order of magnitude (if not, consider changing the moment conditions). This gives consistent estimates of the parameters $\beta$. Then a consistent estimate $\hat{S}$ can be calculated (for instance, with Newey-West).
- Use the consistent $\hat{S}$ from the first step to define a new weighting matrix as $W=$ $\hat{S}^{-1}$. The algorithm is run again to give asymptotically efficient estimates of $\beta$.
- Iterate at least once more. (You may want to consider iterating until the point estimates converge.)

Example 22 (Implementation of 2SLS.) Under the classical 2SLS assumptions, there is no need for iterating since the efficient weighting matrix is $\Sigma_{z z}^{-1} / \sigma^{2}$. Only $\sigma^{2}$ depends on the estimated parameters, but this scaling factor of the loss function does not affect $\hat{\beta}_{2 S L S}$.

One word of warning: if the number of parameters in the covariance matrix $\hat{S}$ is large compared to the number of data points, then $\hat{S}$ tends to be unstable (fluctuates a lot between the steps in the iterations described above) and sometimes also close to singular. The saturation ratio is sometimes used as an indicator of this problem. It is defined as the number of data points of the moment conditions $(q T)$ divided by the number of estimated parameters (the $k$ parameters in $\hat{\beta}$ and the unique $q(q+1) / 2$ parameters in $\hat{S}$ if it is estimated with Newey-West). A value less than 10 is often taken to be an indicator of problems. A possible solution is then to impose restrictions on $S$, for instance, that the autocorrelation is a simple $\operatorname{AR}(1)$ and then estimate $S$ using these restrictions (in which case you cannot use Newey-West, or course).

### 7.8 Testing in GMM

The result in (7.16) can be used to do Wald tests of the parameter vector. For instance, suppose we want to test the $s$ linear restrictions that $R \beta_{0}=r(R$ is $s \times k$ and $r$ is $s \times 1)$ then it must be the case that under null hypothesis

$$
\begin{equation*}
\sqrt{T}(R \hat{\beta}-r) \xrightarrow{d} N\left(\mathbf{0}_{s \times 1}, R V R^{\prime}\right) . \tag{7.23}
\end{equation*}
$$

Remark 23 (Distribution of quadratic forms.) If the $n \times 1$ vector $x \sim N(0, \Sigma)$, then $x^{\prime} \Sigma^{-1} x \sim \chi_{n}^{2}$.

From this remark and the continuous mapping theorem in Remark (19) it follows that, under the null hypothesis that $R \beta_{0}=r$, the Wald test statistics is distributed as a $\chi_{s}^{2}$ variable

$$
\begin{equation*}
T(R \hat{\beta}-r)^{\prime}\left(R V R^{\prime}\right)^{-1}(R \hat{\beta}-r) \xrightarrow{d} \chi_{s}^{2} \tag{7.24}
\end{equation*}
$$

We might also want to test the overidentifying restrictions. The first order conditions (7.11) imply that $k$ linear combinations of the $q$ moment conditions are set to zero by solving for $\hat{\beta}$. Therefore, we have $q-k$ remaining overidentifying restrictions which
should also be close to zero if the model is correct (fits data). Under the null hypothesis that the moment conditions hold (so the overidentifying restrictions hold), we know that $\sqrt{T} \bar{m}\left(\beta_{0}\right)$ is a (scaled) sample average and therefore has (by a CLT) an asymptotic normal distribution. It has a zero mean (the null hypothesis) and the covariance matrix in (7.12). In short,

$$
\begin{equation*}
\sqrt{T} \bar{m}\left(\beta_{0}\right) \xrightarrow{d} N\left(\mathbf{0}_{q \times 1}, S_{0}\right) . \tag{7.25}
\end{equation*}
$$

If would then perhaps be natural to expect that the quadratic form $T \bar{m}(\hat{\beta})^{\prime} S_{0}^{-1} \bar{m}(\hat{\beta})$ should be converge in distribution to a $\chi_{q}^{2}$ variable. That is not correct, however, since $\hat{\beta}$ chosen is such a way that $k$ linear combinations of the first order conditions always (in every sample) are zero. There are, in effect, only $q-k$ nondegenerate random variables in the quadratic form (see Davidson and MacKinnon (1993) 17.6 for a detailed discussion). The correct result is therefore that if we have used optimal weight matrix is used, $W=$ $S_{0}^{-1}$, then

$$
\begin{equation*}
T \bar{m}(\hat{\beta})^{\prime} S_{0}^{-1} \bar{m}(\hat{\beta}) \xrightarrow{d} \chi_{q-k}^{2}, \text { if } W=S_{0}^{-1} \tag{7.26}
\end{equation*}
$$

The left hand side equals $T$ times of value of the loss function (7.8) evaluated at the point estimates, so we could equivalently write what is often called the $J$ test

$$
\begin{equation*}
T J(\hat{\beta}) \sim \chi_{q-k}^{2}, \text { if } W=S_{0}^{-1} \tag{7.27}
\end{equation*}
$$

This also illustrates that with no overidentifying restrictions (as many moment conditions as parameters) there are, of course, no restrictions to test. Indeed, the loss function value is then always zero at the point estimates.

Example 24 (Test of overidentifying assumptions in $2 S L S$.) In contrast to the IV method, 2 SLS allows us to test overidentifying restrictions (we have more moment conditions than parameters, that is, more instruments than regressors). This is a test of whether the residuals are indeed uncorrelated with all the instruments. If not, the model should be rejected. It can be shown that test (7.27) is (asymptotically, at least) the same as the traditional (Sargan (1964), see Davidson (2000) 8.4) test of the overidentifying restrictions in 2SLS. In the latter, the fitted residuals are regressed on the instruments; $T R^{2}$ from that regression is $\chi^{2}$ distributed with as many degrees of freedom as the number of overidentifying restrictions.

Example 25 (Results from GMM on CCAPM; continuing Example 6.) The instruments
could be anything known at $t$ or earlier could be used as instruments. Actually, Hansen and Singleton (1982) and Hansen and Singleton (1983) use lagged $R_{i, t+1} c_{t+1} / c_{t}$ as instruments, and estimate $\gamma$ to be 0.68 to 0.95 , using monthly data. However, $T J_{T}(\hat{\beta})$ is large and the model can usually be rejected at the 5\% significance level. The rejection is most clear when multiple asset returns are used. If T-bills and stocks are tested at the same time, then the rejection would probably be overwhelming.

Another test is to compare a restricted and a less restricted model, where we have used the optimal weighting matrix for the less restricted model in estimating both the less restricted and more restricted model (the weighting matrix is treated as a fixed matrix in the latter case). It can be shown that the test of the $s$ restrictions (the "D test", similar in flavour to an LR test), is

$$
\begin{equation*}
T\left[J\left(\hat{\beta}^{\text {restricted }}\right)-J\left(\hat{\beta}^{\text {less restricted }}\right)\right] \sim \chi_{s}^{2}, \text { if } W=S_{0}^{-1} . \tag{7.28}
\end{equation*}
$$

The weighting matrix is typically based on the unrestricted model. Note that (7.27) is a special case, since the model with allows $q$ non-zero parameters (as many as the moment conditions) always attains $J=0$, and that by imposing $s=q-k$ restrictions we get a restricted model.

### 7.9 GMM with Sub-Optimal Weighting Matrix*

When the optimal weighting matrix is not used, that is, when (7.15) does not hold, then the asymptotic covariance matrix of the parameters is given by (7.22) instead of the result in (7.16). That is,

$$
\begin{equation*}
\sqrt{T}\left(\hat{\beta}-\beta_{0}\right) \xrightarrow{d} N\left(\mathbf{0}_{k \times 1}, V_{2}\right), \text { where } V_{2}=\left(D_{0}^{\prime} W D_{0}\right)^{-1} D_{0}^{\prime} W S_{0} W^{\prime} D_{0}\left(D_{0}^{\prime} W D_{0}\right)^{-1} \tag{7.29}
\end{equation*}
$$

The consistency property is not affected.
The test of the overidentifying restrictions (7.26) and (7.27) are not longer valid. Instead, the result is that

$$
\begin{align*}
\sqrt{T} \bar{m}(\hat{\beta}) & \rightarrow^{d} N\left(\mathbf{0}_{q \times 1}, \Psi\right), \text { with } \\
\Psi & =\left[I-D_{0}\left(D_{0}^{\prime} W D_{0}\right)^{-1} D_{0}^{\prime} W\right] S_{0}\left[I-D_{0}\left(D_{0}^{\prime} W D_{0}\right)^{-1} D_{0}^{\prime} W\right]^{\prime} . \tag{7.31}
\end{align*}
$$

This covariance matrix has rank $q-k$ (the number of overidentifying restriction). This distribution can be used to test hypotheses about the moments, for instance, that a particular moment condition is zero.

Proof. (Sketch of proof of (7.30)-(7.31)) Use (7.19) in (7.17) to get

$$
\begin{aligned}
\sqrt{T} \bar{m}(\hat{\beta}) & =\sqrt{T} \bar{m}\left(\beta_{0}\right)+\sqrt{T} \frac{\partial \bar{m}\left(\beta_{1}\right)}{\partial \beta^{\prime}} \Gamma \bar{m}\left(\beta_{0}\right) \\
& =\left[I+\frac{\partial \bar{m}\left(\beta_{1}\right)}{\partial \beta^{\prime}} \Gamma\right] \sqrt{T} \bar{m}\left(\beta_{0}\right)
\end{aligned}
$$

The term in brackets has a probability limit, which by (7.20) equals $I-D_{0}\left(D_{0}^{\prime} W D_{0}\right)^{-1} D_{0}^{\prime} W$. Since $\sqrt{T} \bar{m}\left(\beta_{0}\right) \rightarrow^{d} N\left(\mathbf{0}_{q \times 1}, S_{0}\right)$ we get (7.30).

Remark 26 If the $n \times 1$ vector $X \sim N(0, \Sigma)$, where $\Sigma$ has rank $r \leq n$ then $Y=$ $X^{\prime} \Sigma^{+} X \sim \chi_{r}^{2}$ where $\Sigma^{+}$is the pseudo inverse of $\Sigma$.

Remark 27 The symmetric $\Sigma$ can be decomposed as $\Sigma=Z \Lambda Z^{\prime}$ where $Z$ are the orthogonal eigenvector $\left(Z^{\prime} Z=I\right)$ and $\Lambda$ is a diagonal matrix with the eigenvalues along the main diagonal. The pseudo inverse can then be calculated as $\Sigma^{+}=Z \Lambda^{+} Z^{\prime}$, where

$$
\Lambda^{+}=\left[\begin{array}{cc}
\Lambda_{11}^{-1} & \mathbf{0} \\
\mathbf{0} & \mathbf{0}
\end{array}\right]
$$

with the reciprocals of the non-zero eigen values along the principal diagonal of $\Lambda_{11}^{-1}$.
This remark and (7.31) implies that the test of overidentifying restrictions (Hansen's $J$ statistics) analogous to (7.26) is

$$
\begin{equation*}
T \bar{m}(\hat{\beta})^{\prime} \Psi^{+} \bar{m}(\hat{\beta}) \xrightarrow{d} \chi_{q-k}^{2} . \tag{7.32}
\end{equation*}
$$

It requires calculation of a generalized inverse (denoted by superscript ${ }^{+}$), but this is fairly straightforward since $\Psi$ is a symmetric matrix. It can be shown (a bit tricky) that this simplifies to (7.26) when the optimal weighting matrix is used.

### 7.10 GMM without a Loss Function*

Suppose we sidestep the whole optimization issue and instead specify $k$ linear combinations (as many as there are parameters) of the $q$ moment conditions directly. That is,
instead of the first order conditions (7.11) we postulate that the estimator should solve

$$
\begin{equation*}
\mathbf{0}_{k \times 1}=\underbrace{A}_{k \times q} \underbrace{\bar{m}(\hat{\beta})}_{q \times 1}(\hat{\beta} \text { is } k \times 1) \tag{7.33}
\end{equation*}
$$

The matrix $A$ is chosen by the researcher and it must have rank $k$ (lower rank means that we effectively have too few moment conditions to estimate the $k$ parameters in $\beta$ ). If $A$ is random, then it should have a finite probability limit $A_{0}$ (also with rank $k$ ). One simple case when this approach makes sense is when we want to use a subset of the moment conditions to estimate the parameters (some columns in $A$ are then filled with zeros), but we want to study the distribution of all the moment conditions.

By comparing (7.11) and (7.33) we see that $A$ plays the same role as $\left[\partial \bar{m}(\hat{\beta}) / \partial \beta^{\prime}\right]^{\prime} W$, but with the difference that $A$ is chosen and not allowed to depend on the parameters. In the asymptotic distribution, it is the probability limit of these matrices that matter, so we can actually substitute $A_{0}$ for $D_{0}^{\prime} W$ in the proof of the asymptotic distribution. The covariance matrix in (7.29) then becomes

$$
\begin{align*}
\mathrm{ACov}\left[\sqrt{T}\left(\hat{\beta}-\beta_{0}\right)\right] & =\left(A_{0} D_{0}\right)^{-1} A_{0} S_{0}\left[\left(A_{0} D_{0}\right)^{-1} A_{0}\right]^{\prime} \\
& =\left(A_{0} D_{0}\right)^{-1} A_{0} S_{0} A_{0}^{\prime}\left[\left(A_{0} D_{0}\right)^{-1}\right]^{\prime}
\end{align*}
$$

which can be used to test hypotheses about the parameters
Similarly, the covariance matrix in (7.30) becomes

$$
\begin{equation*}
\mathrm{ACov}[\sqrt{T} \bar{m}(\hat{\beta})]=\left[I-D_{0}\left(A_{0} D_{0}\right)^{-1} A_{0}\right] S_{0}\left[I-D_{0}\left(A_{0} D_{0}\right)^{-1} A_{0}\right]^{\prime} \tag{7.35}
\end{equation*}
$$

which still has reduced rank. As before, this covariance matrix can be used to construct both $t$ type and $\chi^{2}$ tests of the moment conditions.

### 7.11 Simulated Moments Estimator*

## Reference: Ingram and Lee (1991)

It sometimes happens that it is not possible to calculate the theoretical moments in GMM explicitly. For instance, suppose we want to match the variance of the model with the variance of data

$$
\mathrm{E} f\left(x_{t}, z_{t}, \beta_{0}\right)=\mathrm{E}\left(x_{t}-\mu\right)^{2}-\operatorname{Var} \text { _of_model }\left(\beta_{0}\right)=0
$$

but the model is so non-linear that we cannot find a closed form expression for Var_of_model $\left(\beta_{0}\right)$.
The SME involves ( $i$ ) drawing a set of random numbers for the stochastic shocks in the model; (ii) for a given set of parameter values generate a model simulation with $T_{\text {sin }}$ observations, calculating the moments and using those instead of Var_of_model ( $\beta_{0}$ ) (or similarly for other moments), which is then used to evaluate the loss function $J_{T}$. This is repeated for various sets of parameter values until we find the one which minimizes $J_{T}$.

Basically all GMM results go through, but the covariance matrix should be scaled up with $1+T / T_{\text {sim }}$, where $T$ is the sample length. Note that one must use exactly the same random numbers for all simulations

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## 8 Examples and Applications of GMM

### 8.1 GMM and Classical Econometrics: Examples

### 8.1.1 The LS Estimator (General)

The model is

$$
\begin{equation*}
y_{t}=x_{t}^{\prime} \beta_{0}+u_{t} \tag{8.1}
\end{equation*}
$$

where $\beta$ is a $k \times 1$ vector.
The $k$ moment conditions are

$$
\begin{equation*}
\bar{m}(\beta)=\frac{1}{T} \sum_{t=1}^{T} x_{t}\left(y_{t}-x_{t}^{\prime} \beta\right)=\frac{1}{T} \sum_{t=1}^{T} x_{t} y_{t}-\frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime} \beta . \tag{8.2}
\end{equation*}
$$

The point estimates are found by setting all moment conditions to zero (the model is exactly identified), $\bar{m}(\beta)=\mathbf{0}_{k \times 1}$, which gives

$$
\begin{equation*}
\hat{\beta}=\left(\frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right)^{-1} \frac{1}{T} \sum_{t=1}^{T} x_{t} y_{t} \beta \tag{8.3}
\end{equation*}
$$

If we define

$$
\begin{align*}
S_{0} & =\operatorname{ACov}\left[\sqrt{T} \bar{m}\left(\beta_{0}\right)\right]=\operatorname{ACov}\left(\frac{\sqrt{T}}{T} \sum_{t=1}^{T} x_{t} u_{t}\right)  \tag{8.4}\\
D_{0} & =\operatorname{plim} \frac{\partial \bar{m}\left(\beta_{0}\right)}{\partial \beta^{\prime}}=\operatorname{plim}\left(-\frac{1}{T} \sum_{t=1}^{T} x_{t} x^{\prime}\right)=-\Sigma_{x x} \tag{8.5}
\end{align*}
$$

then the asymptotic covariance matrix of $\sqrt{T}\left(\hat{\beta}-\beta_{0}\right)$

$$
\begin{equation*}
V_{L S}=\left(D_{0}^{\prime} S_{0}^{-1} D_{0}\right)^{-1}=\left(\Sigma_{x x}^{\prime} S_{0}^{-1} \Sigma_{x x}\right)^{-1}=\Sigma_{x x}^{-1} S_{0} \Sigma_{x x}^{-1} \tag{8.6}
\end{equation*}
$$

We can then either try to estimate $S_{0}$ by Newey-West, or make further assumptions to simplify $S_{0}$ (see below).

### 8.1.2 The IV/2SLS Estimator (General)

The model is (8.1), but we use an IV/2SLS method. The $q$ moment conditions (with $q \geq k)$ are

$$
\begin{equation*}
\bar{m}(\beta)=\frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \beta\right)=\frac{1}{T} \sum_{t=1}^{T} z_{t} y_{t}-\frac{1}{T} \sum_{t=1}^{T} z_{t} x_{t}^{\prime} \beta \tag{8.7}
\end{equation*}
$$

The loss function is (for some positive definite weighting matrix $W$, not necessarily the optimal)

$$
\begin{equation*}
\bar{m}(\beta)^{\prime} W \bar{m}(\beta)=\left[\frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \beta\right)\right]^{\prime} W\left[\frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \beta\right)\right] \tag{8.8}
\end{equation*}
$$

and the $k$ first order conditions, $\left(\partial \bar{m}(\hat{\beta}) / \partial \beta^{\prime}\right)^{\prime} W \bar{m}(\hat{\beta})=0$, are

$$
\begin{align*}
\mathbf{0}_{k \times 1} & =\left[\frac{\partial}{\partial \beta^{\prime}} \frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \hat{\beta}\right)\right]^{\prime} W \frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \hat{\beta}\right) \\
& =\left[-\frac{1}{T} \sum_{t=1}^{T} z_{t} x_{t}^{\prime}\right]^{\prime} W \frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{t}-x_{t}^{\prime} \hat{\beta}\right) \\
& =-\hat{\Sigma}_{x z} W\left(\hat{\Sigma}_{z y}-\hat{\Sigma}_{z x} \hat{\beta}\right) \tag{8.9}
\end{align*}
$$

We solve for $\hat{\beta}$ as

$$
\begin{equation*}
\hat{\beta}=\left(\hat{\Sigma}_{x z} W \hat{\Sigma}_{z x}\right)^{-1} \hat{\Sigma}_{x z} W \hat{\Sigma}_{z y} \tag{8.10}
\end{equation*}
$$

Define

$$
\begin{align*}
S_{0} & =\operatorname{ACov}\left[\sqrt{T} \bar{m}\left(\beta_{0}\right)\right]=\operatorname{ACov}\left(\frac{\sqrt{T}}{T} \sum_{t=1}^{T} z_{t} u_{t}\right)  \tag{8.11}\\
D_{0} & =\operatorname{plim} \frac{\partial \bar{m}\left(\beta_{0}\right)}{\partial \beta^{\prime}}=\operatorname{plim}\left(-\frac{1}{T} \sum_{t=1}^{T} z_{t} x_{t}^{\prime}\right)=-\Sigma_{z x} \tag{8.12}
\end{align*}
$$

This gives the asymptotic covariance matrix of $\sqrt{T}\left(\hat{\beta}-\beta_{0}\right)$

$$
\begin{equation*}
V=\left(D_{0}^{\prime} S_{0}^{-1} D_{0}\right)^{-1}=\left(\Sigma_{z x}^{\prime} S_{0}^{-1} \Sigma_{z x}\right)^{-1} \tag{8.13}
\end{equation*}
$$

When the model is exactly identified ( $q=k$ ), then we can make some simplifications since $\hat{\Sigma}_{x z}$ is then invertible. This is the case of the classical IV estimator. We get

$$
\begin{equation*}
\hat{\beta}=\hat{\Sigma}_{z x}^{-1} \hat{\Sigma}_{z y} \text { and } V=\Sigma_{z x}^{-1} S_{0}\left(\Sigma_{z x}^{\prime}\right)^{-1} \text { if } q=k \tag{8.14}
\end{equation*}
$$

(Use the rule $(A B C)^{-1}=C^{-1} B^{-1} A^{-1}$ to show this.)

### 8.1.3 Classical LS Assumptions

Reference: Greene (2000) 9.4 and Hamilton (1994) 8.2.
This section returns to the LS estimator in Section (8.1.1) in order to highlight the classical LS assumptions that give the variance matrix $\sigma^{2} \Sigma_{x x}^{-1}$.

We allow the regressors to be stochastic, but require that $x_{t}$ is independent of all $u_{t+s}$ and that $u_{t}$ is iid. It rules out, for instance, that $u_{t}$ and $x_{t-2}$ are correlated and also that the variance of $u_{t}$ depends on $x_{t}$. Expand the expression for $S_{0}$ as

$$
\begin{align*}
S_{0} & =\mathrm{E}\left(\frac{\sqrt{T}}{T} \sum_{t=1}^{T} x_{t} u_{t}\right)\left(\frac{\sqrt{T}}{T} \sum_{t=1}^{T} u_{t} x_{t}^{\prime}\right)  \tag{8.15}\\
& =\frac{1}{T} \mathrm{E}\left(\ldots+x_{s-1} u_{s-1}+x_{s} u_{s}+\ldots\right)\left(\ldots+u_{s-1} x_{s-1}^{\prime}+u_{s} x_{s}^{\prime}+\ldots\right)
\end{align*}
$$

Note that
$\mathrm{E} x_{t-s} u_{t-s} u_{t} x_{t}^{\prime}=\mathrm{E} x_{t-s} x_{t}^{\prime} \mathrm{E} u_{t-s} u_{t}$ (since $u_{t}$ and $x_{t-s}$ independent)

$$
=\left\{\begin{array}{l}
0 \text { if } s \neq 0\left(\text { since } \mathrm{E} u_{s-1} u_{s}=0 \text { by iid } u_{t}\right)  \tag{8.16}\\
\mathrm{E} x_{t} x_{t}^{\prime} \mathrm{E} u_{t} u_{t} \text { else }
\end{array}\right.
$$

This means that all cross terms (involving different observations) drop out and that we can write

$$
\begin{align*}
S_{0} & =\frac{1}{T} \sum_{t=1}^{T} \mathrm{E} x_{t} x_{t}^{\prime} \mathrm{E} u_{t}^{2}  \tag{8.17}\\
& =\sigma^{2} \frac{1}{T} \mathrm{E} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\left(\text { since } u_{t} \text { is iid and } \sigma^{2}=\mathrm{E} u_{t}^{2}\right)  \tag{8.18}\\
& =\sigma^{2} \Sigma_{x x} \tag{8.19}
\end{align*}
$$

Using this in (8.6) gives

$$
\begin{equation*}
V=\sigma^{2} \Sigma_{x x}^{-1} \tag{8.20}
\end{equation*}
$$

### 8.1.4 Almost Classical LS Assumptions: White's Heteroskedasticity.

Reference: Greene (2000) 12.2 and Davidson and MacKinnon (1993) 16.2.
The only difference compared with the classical LS assumptions is that $u_{t}$ is now allowed to be heteroskedastic, but this heteroskedasticity is not allowed to depend on the moments of $x_{t}$. This means that (8.17) holds, but (8.18) does not since $\mathrm{E} u_{t}^{2}$ is not the same for all $t$.

However, we can still simplify (8.17) a bit more. We assumed that $\mathrm{E} x_{t} x_{t}^{\prime}$ and $\mathrm{E} u_{t}^{2}$ (which can both be time varying) are not related to each other, so we could perhaps multiply $\mathrm{E} x_{t} x_{t}^{\prime}$ by $\Sigma_{t=1}^{T} \mathrm{E} u_{t}^{2} / T$ instead of by $\mathrm{E} u_{t}^{2}$. This is indeed true asymptotically-where any possible "small sample" relation between $\mathrm{E} x_{t} x_{t}^{\prime}$ and $\mathrm{E} u_{t}^{2}$ must wash out due to the assumptions of independence (which are about population moments).

In large samples we therefore have

$$
\begin{align*}
S_{0} & =\left(\frac{1}{T} \sum_{t=1}^{T} \mathrm{E} u_{t}^{2}\right)\left(\frac{1}{T} \sum_{t=1}^{T} \mathrm{E} x_{t} x_{t}^{\prime}\right) \\
& =\left(\frac{1}{T} \sum_{t=1}^{T} \mathrm{E} u_{t}^{2}\right)\left(\mathrm{E} \frac{1}{T} \sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right) \\
& =\omega^{2} \Sigma_{x x} \tag{8.21}
\end{align*}
$$

where $\omega^{2}$ is a scalar. This is very similar to the classical LS case, except that $\omega^{2}$ is the average variance of the residual rather than the constant variance. In practice, the estimator of $\omega^{2}$ is the same as the estimator of $\sigma^{2}$, so we can actually apply the standard LS formulas in this case.

This is the motivation for why White's test for heteroskedasticity makes sense: if the heteroskedasticity is not correlated with the regressors, then the standard LS formula is correct (provided there is no autocorrelation).

### 8.1.5 Estimating the Mean of a Process

Suppose $u_{t}$ is heteroskedastic, but not autocorrelated. In the regression $y_{t}=\alpha+u_{t}$, $x_{t}=z_{t}=1$. This is a special case of the previous example, since $\mathrm{E} u_{t}^{2}$ is certainly unrelated to $\mathrm{E} x_{t} x_{t}^{\prime}=1$ (since it is a constant). Therefore, the LS covariance matrix is the correct variance of the sample mean as an estimator of the mean, even if $u_{t}$ are heteroskedastic (provided there is no autocorrelation).

### 8.1.6 The Classical 2SLS Assumptions*

Reference: Hamilton (1994) 9.2.
The classical 2SLS case assumes that $z_{t}$ is independent of all $u_{t+s}$ and that $u_{t}$ is iid. The covariance matrix of the moment conditions are

$$
\begin{equation*}
S_{0}=\mathrm{E}\left(\frac{1}{\sqrt{T}} \sum_{t=1}^{T} z_{t} u_{t}\right)\left(\frac{1}{\sqrt{T}} \sum_{t=1}^{T} u_{t} z_{t}^{\prime}\right) \tag{8.22}
\end{equation*}
$$

so by following the same steps in (8.16)-(8.19) we get $S_{0}=\sigma^{2} \Sigma_{z z}$. The optimal weighting matrix is therefore $W=\Sigma_{z z}^{-1} / \sigma^{2}$ (or $\left(Z^{\prime} Z / T\right)^{-1} / \sigma^{2}$ in matrix form). We use this result in (8.10) to get

$$
\begin{equation*}
\hat{\beta}_{2 S L S}=\left(\hat{\Sigma}_{x z} \hat{\Sigma}_{z z}^{-1} \hat{\Sigma}_{z x}\right)^{-1} \hat{\Sigma}_{x z} \hat{\Sigma}_{z z}^{-1} \hat{\Sigma}_{z y} \tag{8.23}
\end{equation*}
$$

which is the classical 2SLS estimator.
Since this GMM is efficient (for a given set of moment conditions), we have established that 2SLS uses its given set of instruments in the efficient way-provided the classical 2SLS assumptions are correct. Also, using the weighting matrix in (8.13) gives

$$
\begin{equation*}
V=\left(\Sigma_{x z} \frac{1}{\sigma^{2}} \Sigma_{z z}^{-1} \Sigma_{z x}\right)^{-1} \tag{8.24}
\end{equation*}
$$

### 8.2 Identification of Systems of Simultaneous Equations

Reference: Greene (2000) 16.1-3
This section shows how the GMM moment conditions can be used to understand if the parameters in a system of simultaneous equations are identified or not.

The structural model (form) is

$$
\begin{equation*}
F y_{t}+G z_{t}=u_{t}, \tag{8.25}
\end{equation*}
$$

where $y_{t}$ is a vector of endogenous variables, $z_{t}$ a vector of predetermined (exogenous) variables, $F$ is a square matrix, and $G$ is another matrix. ${ }^{1}$ We can write the $j$ th equation of the structural form (8.25) as

$$
\begin{equation*}
y_{j t}=x_{t}^{\prime} \beta+u_{j t}, \tag{8.26}
\end{equation*}
$$

where $x_{t}$ contains the endogenous and exogenous variables that enter the $j$ th equation with non-zero coefficients, that is, subsets of $y_{t}$ and $z_{t}$.

We want to estimate $\beta$ in (8.26). Least squares is inconsistent if some of the regressors are endogenous variables (in terms of (8.25), this means that the $j$ th row in $F$ contains at least one additional non-zero element apart from coefficient on $y_{j t}$ ). Instead, we use IV/2SLS. By assumption, the structural model summarizes all relevant information for the endogenous variables $y_{t}$. This implies that the only useful instruments are the variables in $z_{t}$. (A valid instrument is uncorrelated with the residuals, but correlated with the regressors.) The moment conditions for the $j$ th equation are then
$\mathrm{E} z_{t}\left(y_{j t}-x_{t}^{\prime} \beta\right)=\mathbf{0}$ with sample moment conditions $\frac{1}{T} \sum_{t=1}^{T} z_{t}\left(y_{j t}-x_{t}^{\prime} \beta\right)=\mathbf{0}$.
If there are as many moment conditions as there are elements in $\beta$, then this equation is exactly identified, so the sample moment conditions can be inverted to give the Instrumental variables (IV) estimator of $\beta$. If there are more moment conditions than elements in $\beta$, then this equation is overidentified and we must devise some method for weighting the different moment conditions. This is the 2SLS method. Finally, when there are fewer moment conditions than elements in $\beta$, then this equation is unidentified, and we cannot hope to estimate the structural parameters of it.

We can partition the vector of regressors in (8.26) as $x_{t}^{\prime}=\left[\tilde{z}_{t}^{\prime}, \tilde{y}_{t}^{\prime}\right]$, where $y_{1 t}$ and $z_{1 t}$ are the subsets of $z_{t}$ and $y_{t}$ respectively, that enter the right hand side of (8.26). Partition $z_{t}$ conformably $z_{t}^{\prime}=\left[\tilde{z}_{t}^{\prime}, z_{t}^{* \prime}\right]$, where $z_{t}^{*}$ are the exogenous variables that do not enter (8.26).

[^4] and $\operatorname{Cov}\left(\varepsilon_{t}\right)=F^{-1} \operatorname{Cov}\left(u_{t}\right)\left(F^{-1}\right)^{\prime}$.

We can then rewrite the moment conditions in (8.27) as

$$
\begin{align*}
\mathrm{E} & {\left[\begin{array}{c}
\tilde{z}_{t} \\
z_{t}^{*}
\end{array}\right]\left(y_{j t}-\left[\begin{array}{c}
\tilde{z}_{t} \\
\tilde{y}_{t}
\end{array}\right]^{\prime} \beta\right)=\mathbf{0} . }  \tag{8.28}\\
y_{j t} & =-G_{j} \tilde{z}_{t}-F_{j} \tilde{y}_{t}+u_{j t} \\
& =x_{t}^{\prime} \beta+u_{j t}, \text { where } x_{t}^{\prime}=\left[\tilde{z}_{t}^{\prime}, \tilde{y}_{t}^{\prime}\right], \tag{8.29}
\end{align*}
$$

This shows that we need at least as many elements in $z_{t}^{*}$ as in $\tilde{y}_{t}$ to have this equations identified, which confirms the old-fashioned rule of thumb: there must be at least as many excluded exogenous variables $\left(z_{t}^{*}\right)$ as included endogenous variables ( $\tilde{y}_{t}$ ) to have the equation identified.

This section has discussed identification of structural parameters when 2SLS/IV, one equation at a time, is used. There are other ways to obtain identification, for instance, by imposing restrictions on the covariance matrix. See, for instance, Greene (2000) 16.1-3 for details.

Example 1 (Supply and Demand. Reference: GR 16, Hamilton 9.1.) Consider the simplest simultaneous equations model for supply and demand on a market. Supply is

$$
q_{t}=\gamma p_{t}+u_{t}^{s}, \gamma>0
$$

and demand is

$$
q_{t}=\beta p_{t}+\alpha A_{t}+u_{t}^{d}, \beta<0
$$

where $A_{t}$ is an observable exogenous demand shock (perhaps income). The only meaningful instrument is $A_{t}$. From the supply equation we then get the moment condition

$$
E A_{t}\left(q_{t}-\gamma p_{t}\right)=0
$$

which gives one equation in one unknown, $\gamma$. The supply equation is therefore exactly identified. In contrast, the demand equation is unidentified, since there is only one (meaningful) moment condition

$$
E A_{t}\left(q_{t}-\beta p_{t}-\alpha A_{t}\right)=0,
$$

but two unknowns ( $\beta$ and $\alpha$ ).

Example 2 (Supply and Demand: overidentification.) If we change the demand equation in Example 1 to

$$
q_{t}=\beta p_{t}+\alpha A_{t}+b B_{t}+u_{t}^{d}, \beta<0 .
$$

There are now two moment conditions for the supply curve (since there are two useful instruments)

$$
E\left[\begin{array}{c}
A_{t}\left(q_{t}-\gamma p_{t}\right) \\
B_{t}\left(q_{t}-\gamma p_{t}\right)
\end{array}\right]=\left[\begin{array}{l}
0 \\
0
\end{array}\right],
$$

but still only one parameter: the supply curve is now overidentified. The demand curve is still underidentified (two instruments and three parameters).

### 8.3 Testing for Autocorrelation

This section discusses how GMM can be used to test if a series is autocorrelated. The analysis focuses on first-order autocorrelation, but it is straightforward to extend it to higher-order autocorrelation.

Consider a scalar random variable $x_{t}$ with a zero mean (it is easy to extend the analysis to allow for a non-zero mean). Consider the moment conditions
$m_{t}(\beta)=\left[\begin{array}{l}x_{t}^{2}-\sigma^{2} \\ x_{t} x_{t-1}-\rho \sigma^{2}\end{array}\right]$, so $\bar{m}(\beta)=\frac{1}{T} \sum_{t=1}^{T}\left[\begin{array}{l}x_{t}^{2}-\sigma^{2} \\ x_{t} x_{t-1}-\rho \sigma^{2}\end{array}\right]$, with $\beta=\left[\begin{array}{l}\sigma^{2} \\ \rho\end{array}\right]$
$\sigma^{2}$ is the variance and $\rho$ the first-order autocorrelation so $\rho \sigma^{2}$ is the first-order autocovariance. We want to test if $\rho=0$. We could proceed along two different routes: estimate $\rho$ and test if it is different from zero or set $\rho$ to zero and then test overidentifying restrictions. We analyze how these two approaches work when the null hypothesis of $\rho=0$ is true.

### 8.3.1 Estimating the Autocorrelation Coefficient

We estimate both $\sigma^{2}$ and $\rho$ by using the moment conditions (8.30) and then test if $\rho=$ 0 . To do that we need to calculate the asymptotic variance of $\hat{\rho}$ (there is little hope of being able to calculate the small sample variance, so we have to settle for the asymptotic variance as an approximation).

We have an exactly identified system so the weight matrix does not matter-we can then proceed as if we had used the optimal weighting matrix (all those results apply).

To find the asymptotic covariance matrix of the parameters estimators, we need the probability limit of the Jacobian of the moments and the covariance matrix of the momentsevaluated at the true parameter values. Let $\bar{m}_{i}\left(\beta_{0}\right)$ denote the $i$ th element of the $\bar{m}(\beta)$ vector-evaluated at the true parameter values. The probability of the Jacobian is

$$
D_{0}=\operatorname{plim}\left[\begin{array}{ll}
\partial \bar{m}_{1}\left(\beta_{0}\right) / \partial \sigma^{2} & \partial \bar{m}_{1}\left(\beta_{0}\right) / \partial \rho \\
\partial \bar{m}_{2}\left(\beta_{0}\right) / \partial \sigma^{2} & \partial \bar{m}_{2}\left(\beta_{0}\right) / \partial \rho
\end{array}\right]=\left[\begin{array}{cc}
-1 & 0 \\
-\rho & -\sigma^{2}
\end{array}\right]=\left[\begin{array}{cc}
-1 & 0 \\
0 & -\sigma^{2}
\end{array}\right]
$$

since $\rho=0$ (the true value). Note that we differentiate with respect to $\sigma^{2}$, not $\sigma$, since we treat $\sigma^{2}$ as a parameter.

The covariance matrix is more complicated. The definition is

$$
S_{0}=\mathrm{E}\left[\frac{\sqrt{T}}{T} \sum_{t=1}^{T} m_{t}\left(\beta_{0}\right)\right]\left[\frac{\sqrt{T}}{T} \sum_{t=1}^{T} m_{t}\left(\beta_{0}\right)\right]^{\prime} .
$$

Assume that there is no autocorrelation in $m_{t}\left(\beta_{0}\right)$. We can then simplify as

$$
S_{0}=\mathrm{E} m_{t}\left(\beta_{0}\right) m_{t}\left(\beta_{0}\right)^{\prime}
$$

This assumption is stronger than assuming that $\rho=0$, but we make it here in order to illustrate the asymptotic distribution. To get anywhere, we assume that $x_{t}$ is iid $N\left(0, \sigma^{2}\right)$. In this case (and with $\rho=0$ imposed) we get

$$
\begin{align*}
S_{0} & =\mathrm{E}\left[\begin{array}{c}
x_{t}^{2}-\sigma^{2} \\
x_{t} x_{t-1}
\end{array}\right]\left[\begin{array}{c}
x_{t}^{2}-\sigma^{2} \\
x_{t} x_{t-1}
\end{array}\right]^{\prime}=\mathrm{E}\left[\begin{array}{cc}
\left(x_{t}^{2}-\sigma^{2}\right)^{2} & \left(x_{t}^{2}-\sigma^{2}\right) x_{t} x_{t-1} \\
\left(x_{t}^{2}-\sigma^{2}\right) x_{t} x_{t-1} & \left(x_{t} x_{t-1}\right)^{2}
\end{array}\right] \\
& =\left[\begin{array}{cc}
\mathrm{E} x_{t}^{4}-2 \sigma^{2} \mathrm{E} x_{t}^{2}+\sigma^{4} & 0 \\
0 & \mathrm{E} x_{t}^{2} x_{t-1}^{2}
\end{array}\right]=\left[\begin{array}{cc}
2 \sigma^{4} & 0 \\
0 & \sigma^{4}
\end{array}\right] \tag{8.32}
\end{align*}
$$

To make the simplification in the second line we use the facts that $\mathrm{E} x_{t}^{4}=3 \sigma^{4}$ if $x_{t} \sim$ $N\left(0, \sigma^{2}\right)$, and that the normality and the iid properties of $x_{t}$ together imply $\mathrm{E} x_{t}^{2} x_{t-1}^{2}=$ $\mathrm{E} x_{t}^{2} \mathrm{E} x_{t-1}^{2}$ and $\mathrm{E} x_{t}^{3} x_{t-1}=\mathrm{E} \sigma^{2} x_{t} x_{t-1}=0$.

By combining (8.31) and (8.32) we get that

$$
\begin{align*}
\operatorname{ACov}\left(\sqrt{T}\left[\begin{array}{l}
\hat{\sigma}^{2} \\
\hat{\rho}
\end{array}\right]\right) & =\left(D_{0}^{\prime} S_{0}^{-1} D_{0}\right)^{-1} \\
& =\left(\left[\begin{array}{cc}
-1 & 0 \\
0 & -\sigma^{2}
\end{array}\right]^{\prime}\left[\begin{array}{cc}
2 \sigma^{4} & 0 \\
0 & \sigma^{4}
\end{array}\right]^{-1}\left[\begin{array}{cc}
-1 & 0 \\
0 & -\sigma^{2}
\end{array}\right]\right)^{-1} \\
& =\left[\begin{array}{cc}
2 \sigma^{4} & 0 \\
0 & 1
\end{array}\right] \tag{8.33}
\end{align*}
$$

This shows the standard expression for the uncertainty of the variance and that the $\sqrt{T} \hat{\rho}$. Since GMM estimators typically have an asymptotic distribution we have $\sqrt{T} \hat{\rho} \rightarrow^{d}$ $N(0,1)$, so we can test the null hypothesis of no first-order autocorrelation by the test statistics

$$
\begin{equation*}
T \hat{\rho}^{2} \sim \chi_{1}^{2} \tag{8.34}
\end{equation*}
$$

This is the same as the Box-Ljung test for first-order autocorrelation.
This analysis shows that we are able to arrive at simple expressions for the sampling uncertainty of the variance and the autocorrelation-provided we are willing to make strong assumptions about the data generating process. In particular, ewe assumed that data was iid $N\left(0, \sigma^{2}\right)$. One of the strong points of GMM is that we could perform similar tests without making strong assumptions-provided we use a correct estimator of the asymptotic covariance matrix $S_{0}$ (for instance, Newey-West).

### 8.3.2 Testing the Overidentifying Restriction of No Autocorrelation*

We can estimate $\sigma^{2}$ alone and then test if both moment condition are satisfied at $\rho=0$. There are several ways of doing that, but the perhaps most straightforward is skip the loss function approach to GMM and instead specify the "first order conditions" directly as

$$
\begin{align*}
0 & =A \bar{m} \\
& =\left[\begin{array}{ll}
1 & 0
\end{array}\right] \frac{1}{T} \sum_{t=1}^{T}\left[\begin{array}{c}
x_{t}^{2}-\sigma^{2} \\
x_{t} x_{t-1}
\end{array}\right] \tag{8.35}
\end{align*}
$$

which sets $\hat{\sigma}^{2}$ equal to the sample variance.

The only parameter in this estimation problem is $\sigma^{2}$, so the matrix of derivatives becomes

$$
D_{0}=\operatorname{plim}\left[\begin{array}{l}
\partial \bar{m}_{1}\left(\beta_{0}\right) / \partial \sigma^{2}  \tag{8.36}\\
\partial \bar{m}_{2}\left(\beta_{0}\right) / \partial \sigma^{2}
\end{array}\right]=\left[\begin{array}{c}
-1 \\
0
\end{array}\right] .
$$

By using this result, the $A$ matrix in (8.36) and the $S_{0}$ matrix in (8.32,) it is straighforward to calculate the asymptotic covariance matrix the moment conditions. In general, we have

$$
\begin{equation*}
\mathrm{ACov}[\sqrt{T} \bar{m}(\hat{\beta})]=\left[I-D_{0}\left(A_{0} D_{0}\right)^{-1} A_{0}\right] S_{0}\left[I-D_{0}\left(A_{0} D_{0}\right)^{-1} A_{0}\right]^{\prime} \tag{8.37}
\end{equation*}
$$

The term in brackets is here (since $A_{0}=A$ since it is a matrix with constants)

$$
\underbrace{\left[\begin{array}{ll}
1 & 0  \tag{8.38}\\
0 & 1
\end{array}\right]}_{I_{2}}-\underbrace{\left[\begin{array}{c}
-1 \\
0
\end{array}\right]}_{D_{0}}(\underbrace{\left[\begin{array}{ll}
1 & 0
\end{array}\right]}_{A_{0}} \underbrace{\left[\begin{array}{c}
-1 \\
0
\end{array}\right]}_{D_{0}})^{-1} \underbrace{\left[\begin{array}{ll}
1 & 0
\end{array}\right]}_{A_{0}}=\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right] .
$$

We therefore get

$$
\operatorname{ACov}[\sqrt{T} \bar{m}(\hat{\beta})]=\left[\begin{array}{ll}
0 & 0  \tag{8.39}\\
0 & 1
\end{array}\right]\left[\begin{array}{cc}
2 \sigma^{4} & 0 \\
0 & \sigma^{4}
\end{array}\right]\left[\begin{array}{ll}
0 & 0 \\
0 & 1
\end{array}\right]^{\prime}=\left[\begin{array}{cc}
0 & 0 \\
0 & \sigma^{4}
\end{array}\right]
$$

Note that the first moment condition has no sampling variance at the estimated parameters, since the choice of $\hat{\sigma}^{2}$ always sets the first moment condition equal to zero.

The test of the overidentifying restriction that the second moment restriction is also zero is

$$
\begin{equation*}
T \bar{m}^{\prime}(\mathrm{ACov}[\sqrt{T} \bar{m}(\hat{\beta})])^{+} \bar{m} \sim \chi_{1}^{2} \tag{8.40}
\end{equation*}
$$

where we have to use a generalized inverse if the covariance matrix is singular (which it is in (8.39)).

In this case, we get the test statistics (note the generalized inverse)

$$
T\left[\begin{array}{l}
0 \\
\Sigma_{t=1}^{T} x_{t} x_{t-1} / T
\end{array}\right]^{\prime}\left[\begin{array}{cc}
0 & 0 \\
0 & 1 / \sigma^{4}
\end{array}\right]\left[\begin{array}{l}
0 \\
\Sigma_{t=1}^{T} x_{t} x_{t-1} / T
\end{array}\right]=T \frac{\left[\Sigma_{t=1}^{T} x_{t} x_{t-1} / T\right]^{2}}{\sigma^{4}}
$$

which is the $T$ times the square of the sample covariance divided by $\sigma^{4}$. A sample correlation, $\hat{\rho}$, would satisfy $\Sigma_{t=1}^{T} x_{t} x_{t-1} / T=\hat{\rho} \hat{\sigma}^{2}$, which we can use to rewrite (8.41) as $T \hat{\rho}^{2} \hat{\sigma}^{4} / \sigma^{4}$. By approximating $\sigma^{4}$ by $\hat{\sigma}^{4}$ we get the same test statistics as in (8.34).

### 8.4 Estimating and Testing a Normal Distribution

### 8.4.1 Estimating the Mean and Variance

This section discusses how the GMM framework can be used to test if a variable is normally distributed. The analysis cold easily be changed in order to test other distributions as well.

Suppose we have a sample of the scalar random variable $x_{t}$ and that we want to test if the series is normally distributed. We analyze the asymptotic distribution under the null hypothesis that $x_{t}$ is $N\left(\mu, \sigma^{2}\right)$.

We specify four moment conditions

$$
m_{t}=\left[\begin{array}{l}
x_{t}-\mu  \tag{8.42}\\
\left(x_{t}-\mu\right)^{2}-\sigma^{2} \\
\left(x_{t}-\mu\right)^{3} \\
\left(x_{t}-\mu\right)^{4}-3 \sigma^{4}
\end{array}\right] \text { so } \bar{m}=\frac{1}{T} \sum_{t=1}^{T}\left[\begin{array}{l}
x_{t}-\mu \\
\left(x_{t}-\mu\right)^{2}-\sigma^{2} \\
\left(x_{t}-\mu\right)^{3} \\
\left(x_{t}-\mu\right)^{4}-3 \sigma^{4}
\end{array}\right]
$$

Note that $\mathrm{E} m_{t}=\mathbf{0}_{4 \times 1}$ if $x_{t}$ is normally distributed.
Let $\bar{m}_{i}\left(\beta_{0}\right)$ denote the $i$ th element of the $\bar{m}(\beta)$ vector-evaluated at the true parameter values. The probability of the Jacobian is

$$
\begin{align*}
D_{0} & =\operatorname{plim}\left[\begin{array}{ll}
\partial \bar{m}_{1}\left(\beta_{0}\right) / \partial \mu & \partial \bar{m}_{1}\left(\beta_{0}\right) / \partial \sigma^{2} \\
\partial \bar{m}_{2}\left(\beta_{0}\right) / \partial \mu & \partial \bar{m}_{2}\left(\beta_{0}\right) / \partial \sigma^{2} \\
\partial \bar{m}_{3}\left(\beta_{0}\right) / \partial \mu & \partial \bar{m}_{3}\left(\beta_{0}\right) / \partial \sigma^{2} \\
\partial \bar{m}_{4}\left(\beta_{0}\right) / \partial \mu & \partial \bar{m}_{4}\left(\beta_{0}\right) / \partial \sigma^{2}
\end{array}\right] \\
& =\operatorname{plim} \frac{1}{T} \sum_{t=1}^{T}\left[\begin{array}{cc}
-1 & 0 \\
-2\left(x_{t}-\mu\right) & -1 \\
-3\left(x_{t}-\mu\right)^{2} & 0 \\
-4\left(x_{t}-\mu\right)^{3} & -6 \sigma^{2}
\end{array}\right]=\left[\begin{array}{cc}
-1 & 0 \\
0 & -1 \\
-3 \sigma^{2} & 0 \\
0 & -6 \sigma^{2}
\end{array}\right] \tag{8.43}
\end{align*}
$$

(Recall that we treat $\sigma^{2}$, not $\sigma$, as a parameter.)
The covariance matrix of the scaled moment conditions (at the true parameter values) is

$$
\begin{equation*}
S_{0}=\mathrm{E}\left[\frac{\sqrt{T}}{T} \sum_{t=1}^{T} m_{t}\left(\beta_{0}\right)\right]\left[\frac{\sqrt{T}}{T} \sum_{t=1}^{T} m_{t}\left(\beta_{0}\right)\right]^{\prime}, \tag{8.44}
\end{equation*}
$$

which can be a very messy expression. Assume that there is no autocorrelation in $m_{t}\left(\beta_{0}\right)$, which would certainly be true if $x_{t}$ is iid. We can then simplify as

$$
\begin{equation*}
S_{0}=\mathrm{E} m_{t}\left(\beta_{0}\right) m_{t}\left(\beta_{0}\right)^{\prime}, \tag{8.45}
\end{equation*}
$$

which is the form we use here for illustration. We therefore have (provided $m_{t}\left(\beta_{0}\right)$ is not autocorrelated)
$S_{0}=\mathrm{E}\left(\left[\begin{array}{c}x_{t}-\mu \\ \left(x_{t}-\mu\right)^{2}-\sigma^{2} \\ \left(x_{t}-\mu\right)^{3} \\ \left(x_{t}-\mu\right)^{4}-3 \sigma^{4}\end{array}\right]\right)\left(\left[\begin{array}{c}x_{t}-\mu \\ \left(x_{t}-\mu\right)^{2}-\sigma^{2} \\ \left(x_{t}-\mu\right)^{3} \\ \left(x_{t}-\mu\right)^{4}-3 \sigma^{4}\end{array}\right]\right)^{\prime}=\left[\begin{array}{cccc}\sigma^{2} & 0 & 3 \sigma^{4} & 0 \\ 0 & 2 \sigma^{4} & 0 & 12 \sigma^{6} \\ 3 \sigma^{4} & 0 & 15 \sigma^{6} & 0 \\ 0 & 12 \sigma^{6} & 0 & 96 \sigma^{8}\end{array}\right]$
(8.46)

It is straightforward to derive this result once we have the information in the following remark.

Remark 3 If $X \sim N\left(\mu, \sigma^{2}\right)$, then the first few moments around the mean of a are $\mathrm{E}(X-$ $\mu)=0, \mathrm{E}(X-\mu)^{2}=\sigma^{2}, \mathrm{E}(X-\mu)^{3}=0$ (all odd moments are zero), $\mathrm{E}(X-\mu)^{4}=3 \sigma^{4}$ $\mathrm{E}(X-\mu)^{6}=15 \sigma^{6}$, and $\mathrm{E}(X-\mu)^{8}=105 \sigma^{8}$

Suppose we use the efficient weighting matrix. The asymptotic covariance matrix of the estimated mean and variance is then $\left(\left(D_{0}^{\prime} S_{0}^{-1} D_{0}\right)^{-1}\right)$

$$
\begin{aligned}
\left(\left[\begin{array}{cc}
-1 & 0 \\
0 & -1 \\
-3 \sigma^{2} & 0 \\
0 & -6 \sigma^{2}
\end{array}\right]^{\prime}\left[\begin{array}{cccc}
\sigma^{2} & 0 & 3 \sigma^{4} & 0 \\
0 & 2 \sigma^{4} & 0 & 12 \sigma^{6} \\
3 \sigma^{4} & 0 & 15 \sigma^{6} & 0 \\
0 & 12 \sigma^{6} & 0 & 96 \sigma^{8}
\end{array}\right]^{-1}\left[\begin{array}{cc}
-1 & 0 \\
0 & -1 \\
-3 \sigma^{2} & 0 \\
0 & -6 \sigma^{2}
\end{array}\right]\right)^{-1} & =\left[\begin{array}{cc}
\frac{1}{\sigma^{2}} & 0 \\
0 & \frac{1}{2 \sigma^{4}}
\end{array}\right]^{-1} \\
& =\left[\begin{array}{cc}
\sigma^{2} & 0 \\
0 & 2 \sigma^{4}
\end{array}\right]
\end{aligned}
$$

This is the same as the result from maximum likelihood estimation which use the sample mean and sample variance as the estimators. The extra moment conditions (overidentifying restrictions) does not produce any more efficient estimators-for the simple reason that the first two moments completely characterizes the normal distribution.

### 8.4.2 Testing Normality*

The payoff from the overidentifying restrictions is that we can test if the series is actually normally distributed. There are several ways of doing that, but the perhaps most straightforward is skip the loss function approach to GMM and instead specify the "first order conditions" directly as

$$
\begin{align*}
0 & =A \bar{m} \\
& =\left[\begin{array}{llll}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0
\end{array}\right] \frac{1}{T} \sum_{t=1}^{T}\left[\begin{array}{l}
x_{t}-\mu \\
\left(x_{t}-\mu\right)^{2}-\sigma^{2} \\
\left(x_{t}-\mu\right)^{3} \\
\left(x_{t}-\mu\right)^{4}-3 \sigma^{4}
\end{array}\right] . \tag{8.48}
\end{align*}
$$

The asymptotic covariance matrix the moment conditions is as in (8.37). In this case, the matrix with brackets is
$\underbrace{\left[\begin{array}{llll}1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0 \\ 0 & 0 & 1 & 0 \\ 0 & 0 & 0 & 1\end{array}\right]}_{I_{4}}-\underbrace{\left[\begin{array}{ccc}-1 & 0 \\ 0 & -1 \\ -3 \sigma^{2} & 0 \\ 0 & -6 \sigma^{2}\end{array}\right]}_{D_{0}} \underbrace{\left[\begin{array}{llll}1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0\end{array}\right]}_{A_{0}} \underbrace{\left[\begin{array}{ccc}-1 & 0 \\ 0 & -1 \\ -3 \sigma^{2} & 0 \\ 0 & -6 \sigma^{2}\end{array}\right]}_{D_{0}})^{-1} \underbrace{\left[\begin{array}{llll}1 & 0 & 0 & 0 \\ 0 & 1 & 0 & 0\end{array}\right]}_{A_{0}}$
$=\left[\begin{array}{cccc}0 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 \\ -3 \sigma^{2} & 0 & 1 & 0 \\ 0 & -6 \sigma^{2} & 0 & 1\end{array}\right]$

We therefore get

$$
\begin{align*}
\operatorname{ACov}[\sqrt{T} \bar{m}(\hat{\beta})] & =\left[\begin{array}{ccccc}
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
-3 \sigma^{2} & 0 & 1 & 0 \\
0 & -6 \sigma^{2} & 0 & 1
\end{array}\right]\left[\begin{array}{cccc}
\sigma^{2} & 0 & 3 \sigma^{4} & 0 \\
0 & 2 \sigma^{4} & 0 & 12 \sigma^{6} \\
3 \sigma^{4} & 0 & 15 \sigma^{6} & 0 \\
0 & 12 \sigma^{6} & 0 & 96 \sigma^{8}
\end{array}\right]\left[\begin{array}{cccc}
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
-3 \sigma^{2} & 0 & 1 & 0 \\
0 & -6 \sigma^{2} & 0 & 1
\end{array}\right]^{\prime} \\
& =\left[\begin{array}{cccc}
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
0 & 0 & 6 \sigma^{6} & 0 \\
0 & 0 & 0 & 24 \sigma^{8}
\end{array}\right] \tag{8.50}
\end{align*}
$$

We now form the test statistics for the overidentifying restrictions as in (8.40). In this case, it is (note the generalized inverse)
$T\left[\begin{array}{l}0 \\ 0 \\ \Sigma_{t=1}^{T}\left(x_{t}-\mu\right)^{3} / T \\ \Sigma_{t=1}^{T}\left[\left(x_{t}-\mu\right)^{4}-3 \sigma^{4}\right] / T\end{array}\right]^{\prime}\left[\begin{array}{ccc}0 & 0 & 0 \\ 0 & 0 \\ 0 & 0 & 0 \\ 0 & 0 & 1 /\left(6 \sigma^{6}\right) \\ 0 & 0 & 0 \\ 1 /\left(24 \sigma^{8}\right)\end{array}\right]\left[\begin{array}{l}0 \\ 0 \\ \Sigma_{t=1}^{T}\left(x_{t}-\mu\right)^{3} / T \\ \Sigma_{t=1}^{T}\left[\left(x_{t}-\mu\right)^{4}-3 \sigma^{4}\right] / T\end{array}\right]$
$=\frac{T}{6} \frac{\left[\Sigma_{t=1}^{T}\left(x_{t}-\mu\right)^{3} / T\right]^{2}}{\sigma^{6}}+\frac{T}{24} \frac{\left\{\Sigma_{t=1}^{T}\left[\left(x_{t}-\mu\right)^{4}-3 \sigma^{4}\right] / T\right\}^{2}}{\sigma^{8}}$.

When we approximate $\sigma$ by $\hat{\sigma}$ then this is the same as the Jarque and Bera test of normality.

The analysis shows (once again) that we can arrive at simple closed form results by making strong assumptions about the data generating process. In particular, we assumed that the moment conditions were serially uncorrelated. The GMM test, with a modified estimator of the covariance matrix $S_{0}$, can typically be much more general.

### 8.5 Testing the Implications of an RBC Model

Reference: Christiano and Eichenbaum (1992)
This section shows how the GMM framework can be used to test if an RBC model fits data.

Christiano and Eichenbaum (1992) try to test if the RBC model predictions correspond are significantly different from correlations and variances of data. The first step is to define
a vector of parameters and some second moments

$$
\begin{equation*}
\Psi=\left[\delta, \ldots, \sigma_{\lambda}, \frac{\sigma_{c^{p}}}{\sigma_{y}}, \ldots, \operatorname{Corr}\left(\frac{y}{n}, n\right)\right] \tag{8.52}
\end{equation*}
$$

and estimate it with GMM using moment conditions. One of the moment condition is that the sample average of the labor share in value added equals the coefficient on labor in a Cobb-Douglas production function, another is that just the definitions of a standard deviation, and so forth.

The distribution of the estimator for $\Psi$ is asymptotically normal. Note that the covariance matrix of the moments is calculated similarly to the Newey-West estimator.

The second step is to note that the RBC model generates second moments as a function $h($.$) of the model parameters \left\{\delta, \ldots, \sigma_{\lambda}\right\}$, which are in $\Psi$, that is, the model generated second moments can be thought of as $h(\Psi)$.

The third step is to test if the non-linear restrictions of the model (the model mapping from parameters to second moments) are satisfied. That is, the restriction that the model second moments are as in data

$$
\begin{equation*}
H(\Psi)=h(\Psi)-\left[\frac{\sigma_{c^{p}}}{\sigma_{y}}, \ldots, \operatorname{Corr}\left(\frac{y}{n}, n\right)\right]=0 \tag{8.53}
\end{equation*}
$$

is tested with a Wald test. (Note that this is much like the $R \beta=0$ constraints in the linear case.) From the delta-method we get

$$
\begin{equation*}
\sqrt{T} H(\hat{\Psi}) \xrightarrow{d} N\left(0, \frac{\partial H}{\partial \Psi^{\prime}} \operatorname{Cov}(\hat{\Psi}) \frac{\partial H^{\prime}}{\partial \Psi}\right) . \tag{8.54}
\end{equation*}
$$

Forming the quadratic form

$$
\begin{equation*}
T H(\hat{\Psi})^{\prime}\left(\frac{\partial H}{\partial \Psi^{\prime}} \operatorname{Cov}(\hat{\Psi}) \frac{\partial H^{\prime}}{\partial \Psi}\right)^{-1} H(\hat{\Psi}) \tag{8.55}
\end{equation*}
$$

will as usual give a $\chi^{2}$ distributed test statistic with as many degrees of freedoms as restrictions (the number of functions in (8.53)).

### 8.6 IV on a System of Equations*

Suppose we have two equations

$$
\begin{aligned}
& y_{1 t}=x_{1 t}^{\prime} \beta_{1}+u_{1 t} \\
& y_{2 t}=x_{2 t}^{\prime} \beta_{2}+u_{2 t}
\end{aligned}
$$

and two sets of instruments, $z_{1 t}$ and $z_{2 t}$ with the same dimensions as $x_{1 t}$ and $x_{2 t}$, respectively. The sample moment conditions are

$$
\bar{m}\left(\beta_{1}, \beta_{2}\right)=\frac{1}{T} \sum_{t=1}^{T}\left[\begin{array}{l}
z_{1 t}\left(y_{1 t}-x_{1 t}^{\prime} \beta_{1}\right) \\
z_{2 t}\left(y_{2 t}-x_{2 t}^{\prime} \beta_{2}\right)
\end{array}\right],
$$

Let $\beta=\left(\beta_{1}^{\prime}, \beta_{2}^{\prime}\right)^{\prime}$. Then

$$
\begin{aligned}
\frac{\partial \bar{m}\left(\beta_{1}, \beta_{2}\right)}{\partial \beta^{\prime}} & =\left[\begin{array}{cc}
\frac{\partial}{\partial \beta_{1}} \frac{1}{T} \sum_{t=1}^{T} z_{1 t}\left(y_{1 t}-x_{1 t}^{\prime} \beta_{1}\right) & \frac{\partial}{\partial \beta_{2}^{\prime}} \frac{1}{T} \sum_{t=1}^{T} z_{1 t}\left(y_{1 t}-x_{1 t}^{\prime} \beta_{1}\right) \\
\frac{\partial}{\partial \beta_{1}^{\prime}} \frac{1}{T} \sum_{t=1}^{T} z_{2 t}\left(y_{2 t}-x_{2 t}^{\prime} \beta_{2}\right) & \frac{\partial}{\partial \beta_{2}^{\prime}} \frac{1}{T} \sum_{t=1}^{T} z_{2 t}\left(y_{2 t}-x_{2 t}^{\prime} \beta_{2}\right)
\end{array}\right] \\
& =\left[\begin{array}{cc}
\frac{1}{T} \sum_{t=1}^{T} z_{1 t} x_{1 t}^{\prime} & \mathbf{0} \\
\mathbf{0} & \frac{1}{T} \sum_{t=1}^{T} z_{2 t} x_{2 t}^{\prime}
\end{array}\right]
\end{aligned}
$$

This is invertible so we can premultiply the first order condition with the inverse of $\left[\partial \bar{m}(\beta) / \partial \beta^{\prime}\right]^{\prime} A$ and get $\bar{m}(\beta)=\mathbf{0}_{k \times 1}$. We can solve this system for $\beta_{1}$ and $\beta_{2}$ as

$$
\begin{aligned}
{\left[\begin{array}{l}
\beta_{1} \\
\beta_{2}
\end{array}\right] } & =\left[\begin{array}{cc}
\frac{1}{T} \sum_{t=1}^{T} z_{1 t} x_{1 t}^{\prime} & \mathbf{0} \\
\mathbf{0} & \frac{1}{T} \sum_{t=1}^{T} z_{2 t} x_{2 t}^{\prime}
\end{array}\right]^{-1}\left[\begin{array}{c}
\frac{1}{T} \sum_{t=1}^{T} z_{1 t} y_{1 t} \\
\frac{1}{T} \sum_{t=1}^{T} z_{2 t} y_{2 t}
\end{array}\right] \\
& =\left[\begin{array}{cc}
\left(\frac{1}{T} \sum_{t=1}^{T} z_{1 t} x_{1 t}^{\prime}\right)^{-1} & \mathbf{0} \\
\mathbf{0} & \left(\frac{1}{T} \sum_{t=1}^{T} z_{2 t} x_{2 t}^{\prime}\right)^{-1}
\end{array}\right]\left[\begin{array}{c}
\frac{1}{T} \sum_{t=1}^{T} z_{1 t} y_{1 t} \\
\frac{1}{T} \sum_{t=1}^{T} z_{2 t} y_{2 t}
\end{array}\right] .
\end{aligned}
$$

This is IV on each equation separately, which follows from having an exactly identified system.

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## 11 Vector Autoregression (VAR)

Reference: Hamilton (1994) 10-11; Greene (2000) 17.5; Johnston and DiNardo (1997) 9.1-9.2 and Appendix 9.2; and Pindyck and Rubinfeld (1997) 9.2 and 13.5.

Let $y_{t}$ be an $n \times 1$ vector of variables. The $\operatorname{VAR}(p)$ is

$$
\begin{equation*}
y_{t}=\mu+A_{1} y_{t-1}+\ldots+A_{p} y_{t-p}+\varepsilon_{t}, \varepsilon_{t} \text { is white noise, } \operatorname{Cov}\left(\varepsilon_{t}\right)=\Omega \tag{11.1}
\end{equation*}
$$

Example $1\left(\operatorname{VAR}(2)\right.$ of $2 \times 1$ vector.) Let $y_{t}=\left[\begin{array}{ll}x_{t} & z_{t}\end{array}\right]^{\prime}$. Then
$\left[\begin{array}{c}x_{t} \\ z_{t}\end{array}\right]=\left[\begin{array}{ll}A_{1,11} & A_{1,12} \\ A_{1,21} & A_{1,22}\end{array}\right]\left[\begin{array}{c}x_{t-1} \\ z_{t-1}\end{array}\right]+\left[\begin{array}{ll}A_{2,11} & A_{2,12} \\ A_{2,21} & A_{2,22}\end{array}\right]\left[\begin{array}{c}x_{t-2} \\ z_{t-2}\end{array}\right]+\left[\begin{array}{c}\varepsilon_{1, t} \\ \varepsilon_{2, t}\end{array}\right]$.
Issues:

- Variable selection
- Lag length
- Estimation
- Purpose: data description (Granger-causality, impulse response, forecast error variance decomposition), forecasting, policy analysis (Lucas critique)?


### 11.1 Canonical Form

$\operatorname{A~} \operatorname{VAR}(p)$ can be rewritten as a $\operatorname{VAR}(1)$. For instance, a $\operatorname{VAR}(2)$ can be written as

$$
\begin{align*}
{\left[\begin{array}{l}
y_{t} \\
y_{t-1}
\end{array}\right] } & =\left[\begin{array}{l}
\mu \\
0
\end{array}\right]+\left[\begin{array}{ll}
A_{1} & A_{2} \\
I & 0
\end{array}\right]\left[\begin{array}{l}
y_{t-1} \\
y_{t-2}
\end{array}\right]+\left[\begin{array}{l}
\varepsilon_{t} \\
0
\end{array}\right] \text { or }  \tag{11.3}\\
y_{t}^{*} & =\mu^{*}+A y_{t-1}^{*}+\varepsilon_{t}^{*} . \tag{11.4}
\end{align*}
$$

Example 2 (Canonical form of a univariate $A R(2)$.)

$$
\left[\begin{array}{l}
y_{t} \\
y_{t-1}
\end{array}\right]=\left[\begin{array}{l}
\mu \\
0
\end{array}\right]+\left[\begin{array}{ll}
a_{1} & a_{2} \\
1 & 0
\end{array}\right]\left[\begin{array}{l}
y_{t-1} \\
y_{t-2}
\end{array}\right]+\left[\begin{array}{l}
\varepsilon_{t} \\
0
\end{array}\right]
$$

Example 3 (Canonical for of $\operatorname{VAR}(2)$ of $2 \times 1$ vector.) Continuing on the previous example, we get

$$
\left[\begin{array}{l}
x_{t} \\
z_{t} \\
x_{t-1} \\
z_{t-1}
\end{array}\right]=\left[\begin{array}{llll}
A_{1,11} & A_{1,11} & A_{2,11} & A_{2,12} \\
A_{1,21} & A_{1,22} & A_{2,21} & A_{2,22} \\
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0
\end{array}\right]\left[\begin{array}{l}
x_{t-1} \\
z_{t-1} \\
x_{t-2} \\
z_{t-2}
\end{array}\right]+\left[\begin{array}{l}
\varepsilon_{1, t} \\
\varepsilon_{2, t} \\
0 \\
0
\end{array}\right]
$$

### 11.2 Moving Average Form and Stability

Consider a $\operatorname{VAR}(1)$, or a $\operatorname{VAR}(1)$ representation of $\operatorname{VAR}(p)$ or an $\operatorname{AR}(p)$

$$
\begin{equation*}
y_{t}^{*}=A y_{t-1}^{*}+\varepsilon_{t}^{*} . \tag{11.5}
\end{equation*}
$$

Solve recursively backwards (substitute for $y_{t-s}^{*}=A y_{t-s-1}^{*}+\varepsilon_{t-s}^{*}, s=1,2, \ldots$ ) to get the vector moving average representation (VMA), or impulse response function

$$
\begin{align*}
y_{t}^{*} & =A\left(A y_{t-2}^{*}+\varepsilon_{t-1}^{*}\right)+\varepsilon_{t}^{*} \\
& =A^{2} y_{t-2}^{*}+A \varepsilon_{t-1}^{*}+\varepsilon_{t}^{*} \\
& =A^{2}\left(A y_{t-3}^{*}+\varepsilon_{t-2}^{*}\right)+A \varepsilon_{t-1}^{*}+\varepsilon_{t}^{*} \\
& =A^{3} y_{t-3}^{*}+A^{2} \varepsilon_{t-2}^{*}+A \varepsilon_{t-1}^{*}+\varepsilon_{t}^{*} \\
& \vdots  \tag{11.6}\\
& =A^{K+1} y_{t-K-1}^{*}+\sum_{s=0}^{K} A^{s} \varepsilon_{t-s}^{*} .
\end{align*}
$$

Remark 4 (Spectral decomposition.) The $n$ eigenvalues ( $\lambda_{i}$ ) and associated eigenvectors $\left(z_{i}\right)$ of the $n \times n$ matrix $A$ satisfies

$$
\left(A-\lambda_{i} I_{n}\right) z_{i}=\mathbf{0}_{n \times 1} .
$$

If the eigenvectors are linearly independent, then

$$
A=Z \Lambda Z^{-1} \text {, where } \Lambda=\left[\begin{array}{cccc}
\lambda_{1} & 0 & \cdots & 0 \\
0 & \lambda_{2} & \cdots & 0 \\
\vdots & \vdots & \cdots & \vdots \\
0 & 0 & \cdots & \lambda_{n}
\end{array}\right] \text { and } Z=\left[\begin{array}{llll}
z_{1} & z_{2} & \cdots & z_{n}
\end{array}\right]
$$

Note that we therefore get

$$
A^{2}=A A=Z \Lambda Z^{-1} Z \Lambda Z^{-1}=Z \Lambda \Lambda Z^{-1}=Z \Lambda^{2} Z^{-1} \Rightarrow A^{q}=Z \Lambda^{q} Z^{-1}
$$

Remark 5 (Modulus of complex number.) If $\lambda=a+$ bi, where $i=\sqrt{-1}$, then $|\lambda|=$ $|a+b i|=\sqrt{a^{2}+b^{2}}$.

We want $\lim _{K \rightarrow \infty} A^{K+1} y_{t-K-1}^{*}=0$ (stable VAR) to get a moving average representation of $y_{t}$ (where the influence of the starting values vanishes asymptotically). We note from the spectral decompositions that $A^{K+1}=Z \Lambda^{K+1} Z^{-1}$, where $Z$ is the matrix of eigenvectors and $\Lambda$ a diagonal matrix with eigenvalues. Clearly, $\lim _{K \rightarrow \infty} A^{K+1} y_{t-K-1}^{*}=$ 0 is satisfied if the eigenvalues of $A$ are all less than one in modulus

Example $6\left(A R(1)\right.$.) For the univariate $A R(1) y_{t}=a y_{t-1}+\varepsilon_{t}$, the characteristic equation is $(a-\lambda) z=0$, which is only satisfied if the eigenvalue is $\lambda=a$. The $A R(1)$ is therefore stable (and stationarity) if $-1<a<1$.

If we have a stable VAR, then (11.6) can be written

$$
\begin{align*}
y_{t}^{*} & =\sum_{s=0}^{\infty} A^{s} \varepsilon_{t-s}^{*}  \tag{11.7}\\
& =\varepsilon_{t}^{*}+A \varepsilon_{t-1}^{*}+A^{2} \varepsilon_{t-2}^{*}+\ldots
\end{align*}
$$

We may pick out the first $n$ equations from (11.7) (to extract the "original" variables from the canonical form) and write them as

$$
\begin{equation*}
y_{t}=\varepsilon_{t}+C_{1} \varepsilon_{t-1}+C_{2} \varepsilon_{t-2}+\ldots \tag{11.8}
\end{equation*}
$$

which is the vector moving average, VMA, form of the VAR.
Example 7 (AR(2), Example (2) continued.) Let $\mu=0$ in 2 and note that the VMA of the canonical form is

$$
\left[\begin{array}{l}
y_{t} \\
y_{t-1}
\end{array}\right]=\left[\begin{array}{l}
\varepsilon_{t} \\
0
\end{array}\right]+\left[\begin{array}{ll}
a_{1} & a_{2} \\
1 & 0
\end{array}\right]\left[\begin{array}{l}
\varepsilon_{t-1} \\
0
\end{array}\right]+\left[\begin{array}{cc}
a_{1}^{2}+a_{2} & a_{1} a_{2} \\
a_{1} & a_{2}
\end{array}\right]\left[\begin{array}{l}
\varepsilon_{t-2} \\
0
\end{array}\right]+\ldots
$$

The MA of $y_{t}$ is therefore

$$
y_{t}=\varepsilon_{t}+a_{1} \varepsilon_{t-1}+\left(a_{1}^{2}+a_{2}\right) \varepsilon_{t-2}+\ldots
$$

Note that

$$
\begin{equation*}
\frac{\partial y_{t}}{\partial \varepsilon_{t-s}^{\prime}}=C_{s} \text { or } \frac{\partial \mathrm{E}_{t} y_{t+s}}{\partial \varepsilon_{t}^{\prime}}=C_{s}, \text { with } C_{0}=I \tag{11.9}
\end{equation*}
$$

so the impulse response function is given by $\left\{I, C_{1}, C_{2}, \ldots\right\}$. Note that it is typically only meaningful to discuss impulse responses to uncorrelated shocks with economic interpretations. The idea behind structural VARs (discussed below) is to impose enough restrictions to achieve this.

Example 8 (Impulse response function for $A R(1)$.) Let $y_{t}=\rho y_{t-1}+\varepsilon_{t}$. The MA representation is $y_{t}=\sum_{s=0}^{t} \rho^{s} \varepsilon_{t-s}$, so $\partial y_{t} / \partial \varepsilon_{t-s}=\partial E_{t} y_{t+s} / \partial \varepsilon_{t}=\rho^{s}$. Stability requires $|\rho|<1$, so the effect of the initial value eventually dies off $\left(\lim _{s \rightarrow \infty} \partial y_{t} / \partial \varepsilon_{t-s}=0\right)$.

Example 9 (Numerical VAR(1) of $2 \times 1$ vector.) Consider the VAR(1)

$$
\left[\begin{array}{l}
x_{t} \\
z_{t}
\end{array}\right]=\left[\begin{array}{ll}
0.5 & 0.2 \\
0.1 & -0.3
\end{array}\right]\left[\begin{array}{l}
x_{t-1} \\
z_{t-1}
\end{array}\right]+\left[\begin{array}{l}
\varepsilon_{1, t} \\
\varepsilon_{2, t}
\end{array}\right]
$$

The eigenvalues are approximately 0.52 and -0.32 , so this is a stable VAR. The VMA is

$$
\left[\begin{array}{l}
x_{t} \\
z_{t}
\end{array}\right]=\left[\begin{array}{l}
\varepsilon_{1, t} \\
\varepsilon_{2, t}
\end{array}\right]+\left[\begin{array}{ll}
0.5 & 0.2 \\
0.1 & -0.3
\end{array}\right]\left[\begin{array}{l}
\varepsilon_{1, t-1} \\
\varepsilon_{2, t-1}
\end{array}\right]+\left[\begin{array}{ll}
0.27 & 0.04 \\
0.02 & 0.11
\end{array}\right]\left[\begin{array}{l}
\varepsilon_{1, t-2} \\
\varepsilon_{2, t-2}
\end{array}\right]+\ldots
$$

### 11.3 Estimation

The MLE, conditional on the initial observations, of the VAR is the same as OLS estimates of each equation separately. The MLE of the $i j^{t h}$ element in $\operatorname{Cov}\left(\varepsilon_{t}\right)$ is given by $\sum_{t=1}^{T} \hat{v}_{i t} \hat{v}_{j t} / T$, where $\hat{v}_{i t}$ and $\hat{v}_{j t}$ are the OLS residuals.

Note that the VAR system is a system of "seemingly unrelated regressions," with the same regressors in each equation. The OLS on each equation is therefore the GLS, which coincides with MLE if the errors are normally distributed.

### 11.4 Granger Causality

Main message: Granger-causality might be useful, but it is not the same as causality.
Definition: if $z$ cannot help forecast $x$, then $z$ does not Granger-cause $x$; the MSE of the forecast $\mathrm{E}\left(x_{t} \mid x_{t-s}, z_{t-s}, s>0\right)$ equals the MSE of the forecast $\mathrm{E}\left(x_{t} \mid x_{t-s}, s>0\right)$.

Test: Redefine the dimensions of $x_{t}$ and $z_{t}$ in (11.2): let $x_{t}$ be $n_{1} \times 1$ and $z_{t}$ is $n_{2} \times 1$. If the $n_{1} \times n_{2}$ matrices $A_{1,12}=\mathbf{0}$ and $A_{2,12}=\mathbf{0}$, then $z$ fail to Granger-cause $x$. (In general, we would require $A_{s, 12}=\mathbf{0}$ for $s=1, \ldots, p$.) This carries over to the MA representation in (11.8), so $C_{s, 12}=\mathbf{0}$.

These restrictions can be tested with an F-test. The easiest case is when $x$ is a scalar, since we then simply have a set of linear restrictions on a single OLS regression.

Example 10 (RBC and nominal neutrality.) Suppose we have an RBC model which says that money has no effect on the real variables (for instance, output, capital stock, and the productivity level). Money stock should not Granger-cause real variables.

Example 11 (Granger causality and causality.) Do Christmas cards cause Christmas?
Example 12 (Granger causality and causality II, from Hamilton 11.) Consider the price $P_{t}$ of an asset paying dividends $D_{t}$. Suppose the expected return $\left(E_{t}\left(P_{t+1}+D_{t+1}\right) / P_{t}\right)$ is a constant, R. The price then satisfies $P_{t}=E_{t} \sum_{s=1}^{\infty} R^{-s} D_{t+s}$. Suppose $D_{t}=u_{t}+$ $\delta u_{t-1}+v_{t}$, so $E_{t} D_{t+1}=\delta u_{t}$ and $E_{t} D_{t+s}=0$ for $s>1$. This gives $P_{t}=\delta u_{t} / R$, and $D_{t}=u_{t}+v_{t}+R P_{t-1}$, so the VAR is

$$
\left[\begin{array}{c}
P_{t} \\
D_{t}
\end{array}\right]=\left[\begin{array}{ll}
0 & 0 \\
R & 0
\end{array}\right]\left[\begin{array}{c}
P_{t-1} \\
D_{t-1}
\end{array}\right]+\left[\begin{array}{l}
\delta u_{t} / R \\
u_{t}+v_{t}
\end{array}\right]
$$

where P Granger-causes D. Of course, the true causality is from D to P. Problem: forward looking behavior.

Example 13 (Money and output, Sims (1972).) Sims found that output, y does not Grangercause money, $m$, but that $m$ Granger causes $y$. His interpretation was that money supply is exogenous (set by the Fed) and that money has real effects. Notice how he used a combination of two Granger causality test to make an economic interpretation.

Example 14 (Granger causality and omitted information.*) Consider the VAR

$$
\left[\begin{array}{l}
y_{1 t} \\
y_{2 t} \\
y_{3 t}
\end{array}\right]=\left[\begin{array}{lll}
a_{11} & a_{12} & 0 \\
0 & a_{22} & 0 \\
0 & a_{32} & a_{33}
\end{array}\right]\left[\begin{array}{l}
y_{1 t-1} \\
y_{2 t-1} \\
y_{3 t-1}
\end{array}\right]+\left[\begin{array}{l}
\varepsilon_{1 t} \\
\varepsilon_{2 t} \\
\varepsilon_{3 t}
\end{array}\right]
$$

Notice that $y_{2 t}$ and $y_{3 t}$ do not depend on $y_{1 t-1}$, so the latter should not be able to Grangercause $y_{3 t}$. However, suppose we forget to use $y_{2 t}$ in the regression and then ask if $y_{1 t}$

Granger causes $y_{3 t}$. The answer might very well be yes since $y_{1 t-1}$ contains information about $y_{2 t-1}$ which does affect $y_{3 t}$. (If you let $y_{1 t}$ be money, $y_{2 t}$ be the (autocorrelated) Solow residual, and $y_{3 t}$ be output, then this is a short version of the comment in King (1986) comment on Bernanke (1986) (see below) on why money may appear to Grangercause output). Also note that adding a nominal interest rate to Sims (see above) moneyoutput VAR showed that money cannot be taken to be exogenous.

### 11.5 Forecasts Forecast Error Variance

The error forecast of the $s$ period ahead forecast is

$$
\begin{equation*}
y_{t+s}-\mathrm{E}_{t} y_{t+s}=\varepsilon_{t+s}+C_{1} \varepsilon_{t+s-1}+\ldots+C_{s-1} \varepsilon_{t+1} \tag{11.10}
\end{equation*}
$$

so the covariance matrix of the ( $s$ periods ahead) forecasting errors is

$$
\begin{equation*}
\mathrm{E}\left(y_{t+s}-\mathrm{E}_{t} y_{t+s}\right)\left(y_{t+s}-\mathrm{E}_{t} y_{t+s}\right)^{\prime}=\Omega+C_{1} \Omega C_{1}^{\prime}+\ldots+C_{s-1} \Omega C_{s-1}^{\prime} . \tag{11.11}
\end{equation*}
$$

For a $\operatorname{VAR}(1), C_{s}=A^{s}$, so we have

$$
\begin{equation*}
y_{t+s}-\mathrm{E}_{t} y_{t+s}=\varepsilon_{t+s}+A \varepsilon_{t+s-1}+\ldots+A^{s} \varepsilon_{t+1}, \text { and } \tag{11.12}
\end{equation*}
$$

$\mathrm{E}\left(y_{t+s}-\mathrm{E}_{t} y_{t+s}\right)\left(y_{t+s}-\mathrm{E}_{t} y_{t+s}\right)^{\prime}=\Omega+A \Omega A^{\prime}+\ldots+A^{s-1} \Omega\left(A^{s-1}\right)^{\prime}$.
Note that $\lim _{s \rightarrow \infty} \mathrm{E}_{t} y_{t+s}=0$, that is, the forecast goes to the unconditional mean (which is zero here, since there are no constants - you could think of $y_{t}$ as a deviation from the mean). Consequently, the forecast error becomes the VMA representation (11.8). Similarly, the forecast error variance goes to the unconditional variance.

Example 15 (Unconditional variance of $\operatorname{VAR(1).)~Letting~} s \rightarrow \infty$ in (11.13) gives

$$
\begin{aligned}
E y_{t} y_{t}^{\prime} & =\sum_{s=0}^{\infty} A^{s} \Omega\left(A^{s}\right)^{\prime} \\
& =\Omega+\left[A \Omega A^{\prime}+A^{2} \Omega\left(A^{2}\right)^{\prime}+\ldots\right] \\
& =\Omega+A\left(\Omega+A \Omega A^{\prime}+\ldots\right) A^{\prime} \\
& =\Omega+A\left(E y_{t} y_{t}^{\prime}\right) A^{\prime}
\end{aligned}
$$

which suggests that we can calculate $E y_{t} y_{t}^{\prime}$ by an iteration (backwards in time) $\Phi_{t}=$ $\Omega+A \Phi_{t+1} A^{\prime}$, starting from $\Phi_{T}=I$, until convergence.

### 11.6 Forecast Error Variance Decompositions*

If the shocks are uncorrelated, then it is often useful to calculate the fraction of $\operatorname{Var}\left(y_{i, t+s}-\mathrm{E}_{t} y_{i, t+s}\right)$ due to the $j^{t h}$ shock, the forecast error variance decomposition. Suppose the covariance matrix of the shocks, here $\Omega$, is a diagonal $n \times n$ matrix with the variances $\omega_{i i}$ along the diagonal. Let $c_{q i}$ be the $\mathrm{i}^{\text {th }}$ column of $C_{q}$. We then have

$$
\begin{equation*}
C_{q} \Omega C_{q}^{\prime}=\sum_{i=1}^{n} \omega_{i i} c_{q i}\left(c_{q i}\right)^{\prime} . \tag{11.14}
\end{equation*}
$$

Example 16 (Illustration of (11.14) with $n=2$.) Suppose

$$
C_{q}=\left[\begin{array}{ll}
c_{11} & c_{12} \\
c_{21} & c_{22}
\end{array}\right] \text { and } \Omega=\left[\begin{array}{cc}
\omega_{11} & 0 \\
0 & \omega_{22}
\end{array}\right]
$$

then

$$
C_{q} \Omega C_{q}^{\prime}=\left[\begin{array}{cc}
\omega_{11} c_{11}^{2}+\omega_{22} c_{12}^{2} & \omega_{11} c_{11} c_{21}+\omega_{22} c_{12} c_{22} \\
\omega_{11} c_{11} c_{21}+\omega_{22} c_{12} c_{22} & \omega_{11} c_{21}^{2}+\omega_{22} c_{22}^{2}
\end{array}\right]
$$

which should be compared with

$$
\begin{aligned}
& \omega_{11}\left[\begin{array}{l}
c_{11} \\
c_{21}
\end{array}\right]\left[\begin{array}{l}
c_{11} \\
c_{21}
\end{array}\right]^{\prime}+\omega_{22}\left[\begin{array}{l}
c_{12} \\
c_{22}
\end{array}\right]\left[\begin{array}{l}
c_{12} \\
c_{22}
\end{array}\right]^{\prime} \\
& =\omega_{11}\left[\begin{array}{cc}
c_{11}^{2} & c_{11} c_{21} \\
c_{11} c_{21} & c_{21}^{2}
\end{array}\right]+\omega_{22}\left[\begin{array}{cc}
c_{12}^{2} & c_{12} c_{22} \\
c_{12} c_{22} & c_{22}^{2}
\end{array}\right] .
\end{aligned}
$$

Applying this on (11.11) gives

$$
\begin{align*}
\mathrm{E}\left(y_{t+s}-\mathrm{E}_{t} y_{t+s}\right)\left(y_{t+s}-\mathrm{E}_{t} y_{t+s}\right)^{\prime} & =\sum_{i=1}^{n} \omega_{i i} I+\sum_{i=1}^{n} \omega_{i i} c_{1 i}\left(c_{1 i}\right)^{\prime}+\ldots+\sum_{i=1}^{n} \omega_{i i} c_{s-1 i}\left(c_{s-1 i}\right)^{\prime} \\
& =\sum_{i=1}^{n} \omega_{i i}\left[I+c_{1 i}\left(c_{1 i}\right)^{\prime}+\ldots+c_{s-1 i}\left(c_{s-1 i}\right)^{\prime}\right] \tag{11.15}
\end{align*}
$$

which shows how the covariance matrix for the $s$-period forecast errors can be decomposed into its $n$ components.

### 11.7 Structural VARs

### 11.7.1 Structural and Reduced Forms

We are usually not interested in the impulse response function (11.8) or the variance decomposition (11.11) with respect to $\varepsilon_{t}$, but with respect to some structural shocks, $u_{t}$, which have clearer interpretations (technology, monetary policy shock, etc.).

Suppose the structural form of the model is

$$
\begin{equation*}
F y_{t}=\alpha+B_{1} y_{t-1}+\ldots+B_{p} y_{t-p}+u_{t}, u_{t} \text { is white noise, } \operatorname{Cov}\left(u_{t}\right)=D \tag{11.16}
\end{equation*}
$$

This could, for instance, be an economic model derived from theory.
Provided $F^{-1}$ exists, it is possible to write the time series process as

$$
\begin{align*}
y_{t} & =F^{-1} \alpha+F^{-1} B_{1} y_{t-1}+\ldots+F^{-1} B_{p} y_{t-p}+F^{-1} u_{t}  \tag{11.17}\\
& =\mu+A_{1} y_{t-1}+\ldots+A_{p} y_{t-p}+\varepsilon_{t}, \operatorname{Cov}\left(\varepsilon_{t}\right)=\Omega \tag{11.18}
\end{align*}
$$

where

$$
\begin{equation*}
\mu=F^{-1} \alpha, A_{s}=F^{-1} B_{s}, \text { and } \varepsilon_{t}=F^{-1} u_{t} \text { so } \Omega=F^{-1} D\left(F^{-1}\right)^{\prime} \tag{11.19}
\end{equation*}
$$

Equation (11.18) is a VAR model, so a VAR can be thought of as a reduced form of the structural model (11.16).

The key to understanding the relation between the structural model and the VAR is the $F$ matrix, which controls how the endogenous variables, $y_{t}$, are linked to each other contemporaneously. In fact, identification of a VAR amounts to choosing an $F$ matrix. Once that is done, impulse responses and forecast error variance decompositions can be made with respect to the structural shocks. For instance, the impulse response function of

[^5]the VAR, (11.8), can be rewritten in terms of $u_{t}=F \varepsilon_{t}($ from (11.19))
\[

$$
\begin{align*}
y_{t} & =\varepsilon_{t}+C_{1} \varepsilon_{t-1}+C_{2} \varepsilon_{t-2}+\ldots \\
& =F^{-1} F \varepsilon_{t}+C_{1} F^{-1} F \varepsilon_{t-1}+C_{2} F^{-1} F \varepsilon_{t-2}+\ldots \\
& =F^{-1} u_{t}+C_{1} F^{-1} u_{t-1}+C_{2} F^{-1} u_{t-2}+\ldots \tag{11.20}
\end{align*}
$$
\]

Remark 17 The easiest way to calculate this representation is by first finding $F^{-1}$ (see below), then writing (11.18) as

$$
\begin{equation*}
y_{t}=\mu+A_{1} y_{t-1}+\ldots+A_{p} y_{t-p}+F^{-1} u_{t} . \tag{11.21}
\end{equation*}
$$

To calculate the impulse responses to the first element in $u_{t}$, set $y_{t-1}, \ldots, y_{t-p}$ equal to the long-run average, $\left(I-A_{1}-\ldots-A p\right)^{-1} \mu$, make the first element in $u_{t}$ unity and all other elements zero. Calculate the response by iterating forward on (11.21), but putting all elements in $u_{t+1}, u_{t+2}, \ldots$ to zero. This procedure can be repeated for the other elements of $u_{t}$.

We would typically pick $F$ such that the elements in $u_{t}$ are uncorrelated with each other, so they have a clear interpretation.

The VAR form can be estimated directly from data. Is it then possible to recover the structural parameters in (11.16) from the estimated VAR (11.18)? Not without restrictions on the structural parameters in $F, B_{s}, \alpha$, and $D$. To see why, note that in the structural form (11.16) we have $(p+1) n^{2}$ parameters in $\left\{F, B_{1}, \ldots, B_{p}\right\}, n$ parameters in $\alpha$, and $n(n+1) / 2$ unique parameters in $D$ (it is symmetric). In the VAR (11.18) we have fewer parameters: $p n^{2}$ in $\left\{A_{1}, \ldots, A_{p}\right\}, n$ parameters in in $\mu$, and $n(n+1) / 2$ unique parameters in $\Omega$. This means that we have to impose at least $n^{2}$ restrictions on the structural parameters $\left\{F, B_{1}, \ldots, B_{p}, \alpha, D\right\}$ to identify all of them. This means, of course, that many different structural models have can have exactly the same reduced form.

Example 18 (Structural form of the $2 \times 1$ case.) Suppose the structural form of the previous example is
$\left[\begin{array}{ll}F_{11} & F_{12} \\ F_{21} & F_{22}\end{array}\right]\left[\begin{array}{l}x_{t} \\ z_{t}\end{array}\right]=\left[\begin{array}{ll}B_{1,11} & B_{1,12} \\ B_{1,21} & B_{1,22}\end{array}\right]\left[\begin{array}{l}x_{t-1} \\ z_{t-1}\end{array}\right]+\left[\begin{array}{ll}B_{2,11} & B_{2,12} \\ B_{2,21} & B_{2,22}\end{array}\right]\left[\begin{array}{l}x_{t-2} \\ z_{t-2}\end{array}\right]+\left[\begin{array}{l}u_{1, t} \\ u_{2, t}\end{array}\right]$.
This structural form has $3 \times 4+3$ unique parameters. The VAR in (11.2) has $2 \times 4+3$. We need at least 4 restrictions on $\left\{F, B_{1}, B_{2}, D\right\}$ to identify them from $\left\{A_{1}, A_{2}, \Omega\right\}$.

### 11.7.2 "Triangular" Identification 1: Triangular $F$ with $F_{i i}=1$ and Diagonal $D$

## Reference: Sims (1980).

The perhaps most common way to achieve identification of the structural parameters is to restrict the contemporaneous response of the different endogenous variables, $y_{t}$, to the different structural shocks, $u_{t}$. Within in this class of restrictions, the triangular identification is the most popular: assume that $F$ is lower triangular $(n(n+1) / 2$ restrictions) with diagonal element equal to unity, and that $D$ is diagonal ( $n(n-1) / 2$ restrictions), which gives $n^{2}$ restrictions (exact identification).

A lower triangular $F$ matrix is very restrictive. It means that the first variable can react to lags and the first shock, the second variable to lags and the first two shocks, etc. This is a recursive simultaneous equations model, and we obviously need to be careful with how we order the variables. The assumptions that $F_{i i}=1$ is just a normalization.

A diagonal $D$ matrix seems to be something that we would often like to have in a structural form in order to interpret the shocks as, for instance, demand and supply shocks. The diagonal elements of $D$ are the variances of the structural shocks.

Example 19 (Lower triangular F: going from structural form to VAR.) Suppose the structural form is

$$
\left[\begin{array}{cc}
1 & 0 \\
-\alpha & 1
\end{array}\right]\left[\begin{array}{l}
x_{t} \\
z_{t}
\end{array}\right]=\left[\begin{array}{ll}
B_{11} & B_{12} \\
B_{21} & B_{22}
\end{array}\right]\left[\begin{array}{l}
x_{t-1} \\
z_{t-1}
\end{array}\right]+\left[\begin{array}{l}
u_{1, t} \\
u_{2, t}
\end{array}\right]
$$

This is a recursive system where $x_{t}$ does not not depend on the contemporaneous $z_{t}$, and therefore not on the contemporaneous $u_{2 t}$ (see first equation). However, $z_{t}$ does depend on $x_{t}$ (second equation). The VAR (reduced form) is obtained by premultiplying by $F^{-1}$

$$
\begin{aligned}
{\left[\begin{array}{l}
x_{t} \\
z_{t}
\end{array}\right] } & =\left[\begin{array}{ll}
1 & 0 \\
\alpha & 1
\end{array}\right]\left[\begin{array}{ll}
B_{11} & B_{12} \\
B_{21} & B_{22}
\end{array}\right]\left[\begin{array}{l}
x_{t-1} \\
z_{t-1}
\end{array}\right]+\left[\begin{array}{ll}
1 & 0 \\
\alpha & 1
\end{array}\right]\left[\begin{array}{l}
u_{1, t} \\
u_{2, t}
\end{array}\right] \\
& =\left[\begin{array}{ll}
A_{11} & A_{12} \\
A_{21} & A_{22}
\end{array}\right]\left[\begin{array}{l}
x_{t-1} \\
z_{t-1}
\end{array}\right]+\left[\begin{array}{l}
\varepsilon_{1, t} \\
\varepsilon_{2, t}
\end{array}\right] .
\end{aligned}
$$

This means that $\varepsilon_{1 t}=u_{1 t}$, so the first VAR shock equals the first structural shock. In contrast, $\varepsilon_{2, t}=\alpha u_{1, t}+u_{2, t}$, so the second VAR shock is a linear combination of the first
two shocks. The covariance matrix of the VAR shocks is therefore

$$
\operatorname{Cov}\left[\begin{array}{l}
\varepsilon_{1, t} \\
\varepsilon_{2, t}
\end{array}\right]=\left[\begin{array}{ll}
\operatorname{Var}\left(u_{1 t}\right) & \alpha \operatorname{Var}\left(u_{1 t}\right) \\
\alpha \operatorname{Var}\left(u_{1 t}\right) & \alpha^{2} \operatorname{Var}\left(u_{1 t}\right)+\operatorname{Var}\left(u_{2 t}\right)
\end{array}\right] .
$$

This set of identifying restrictions can be implemented by estimating the structural form with LS-equation by equation. The reason is that this is just the old fashioned fully recursive system of simultaneous equations. See, for instance, Greene (2000) 16.3.

### 11.7.3 "Triangular" Identification 2: Triangular $F$ and $D=I$

The identifying restrictions in Section 11.7.2 is actually the same as assuming that $F$ is triangular and that $D=I$. In this latter case, the restriction on the diagonal elements of $F$ has been moved to the diagonal elements of $D$. This is just a change of normalization (that the structural shocks have unit variance). It happens that this alternative normalization is fairly convenient when we want to estimate the VAR first and then recover the structural parameters from the VAR estimates.

Example 20 (Change of normalization in Example 19) Suppose the structural shocks in Example 19 have the covariance matrix

$$
D=\operatorname{Cov}\left[\begin{array}{l}
u_{1, t} \\
u_{2, t}
\end{array}\right]=\left[\begin{array}{ll}
\sigma_{1}^{2} & 0 \\
0 & \sigma_{2}^{2}
\end{array}\right] .
$$

Premultiply the structural form in Example 19 by

$$
\left[\begin{array}{ll}
1 / \sigma_{1} & 0 \\
0 & 1 / \sigma_{2}
\end{array}\right]
$$

to get

$$
\left[\begin{array}{cc}
1 / \sigma_{1} & 0 \\
-\alpha / \sigma_{2} & 1 / \sigma_{2}
\end{array}\right]\left[\begin{array}{l}
x_{t} \\
z_{t}
\end{array}\right]=\left[\begin{array}{ll}
B_{11} / \sigma_{1} & B_{12} / \sigma_{1} \\
B_{21} / \sigma_{2} & B_{22} / \sigma_{2}
\end{array}\right]\left[\begin{array}{l}
x_{t-1} \\
z_{t-1}
\end{array}\right]+\left[\begin{array}{l}
u_{1, t} / \sigma_{1} \\
u_{2, t} / \sigma_{2}
\end{array}\right]
$$

This structural form has a triangular $F$ matrix (with diagonal elements that can be different from unity), and a covariance matrix equal to an identity matrix.

The reason why this alternative normalization is convenient is that it allows us to use the widely available Cholesky decomposition.

Remark 21 (Cholesky decomposition) Let $\Omega$ be an $n \times n$ symmetric positive definite matrix. The Cholesky decomposition gives the unique lower triangular $P$ such that $\Omega=$ $P P^{\prime}$ (some software returns an upper triangular matrix, that is, $Q$ in $\Omega=Q^{\prime} Q$ instead).

Remark 22 Note the following two important features of the Cholesky decomposition. First, each column of $P$ is only identified up to a sign transformation; they can be reversed at will. Second, the diagonal elements in $P$ are typically not unity

Remark 23 (Changing sign of column and inverting.) Suppose the square matrix $A_{2}$ is the same as $A_{1}$ except that the $i^{\text {th }}$ and $j^{\text {th }}$ columns have the reverse signs. Then $A_{2}^{-1}$ is the same as $A_{1}^{-1}$ except that the $i^{\text {th }}$ and $j^{\text {th }}$ rows have the reverse sign

This set of identifying restrictions can be implemented by estimating the VAR with LS and then take the following steps.

- Step 1. From (11.19) $\Omega=F^{-1} I\left(F^{-1}\right)^{\prime}$ (recall $D=I$ is assumed), so a Cholesky decomposition recovers $F^{-1}$ (lower triangular $F$ gives a similar structure of $F^{-1}$, and vice versa, so this works). The signs of each column of $F^{-1}$ can be chosen freely, for instance, so that a productivity shock gets a positive, rather than negative, effect on output. Invert $F^{-1}$ to get $F$.
- Step 2. Invert the expressions in (11.19) to calculate the structural parameters from the VAR parameters as $\alpha=F \mu$, and $B_{s}=F A_{s}$.

Example 24 (Identification of the $2 \times 1$ case.) Suppose the structural form of the previous example is
$\left[\begin{array}{ll}F_{11} & 0 \\ F_{21} & F_{22}\end{array}\right]\left[\begin{array}{l}x_{t} \\ z_{t}\end{array}\right]=\left[\begin{array}{ll}B_{1,11} & B_{1,12} \\ B_{1,21} & B_{1,22}\end{array}\right]\left[\begin{array}{l}x_{t-1} \\ z_{t-1}\end{array}\right]+\left[\begin{array}{ll}B_{2,11} & B_{2,12} \\ B_{2,21} & B_{2,22}\end{array}\right]\left[\begin{array}{l}x_{t-2} \\ z_{t-2}\end{array}\right]+\left[\begin{array}{l}u_{1, t} \\ u_{2, t}\end{array}\right]$, with $D=\left[\begin{array}{ll}1 & 0 \\ 0 & 1\end{array}\right]$.

Step 1 above solves

$$
\begin{aligned}
{\left[\begin{array}{ll}
\Omega_{11} & \Omega_{12} \\
\Omega_{12} & \Omega_{22}
\end{array}\right] } & =\left[\begin{array}{ll}
F_{11} & 0 \\
F_{21} & F_{22}
\end{array}\right]^{-1}\left(\left[\begin{array}{ll}
F_{11} & 0 \\
F_{21} & F_{22}
\end{array}\right]^{-1}\right)^{\prime} \\
& =\left[\begin{array}{cc}
\frac{1}{F_{11}^{2}} & -\frac{F_{21}}{F_{11}^{2} F_{22}} \\
-\frac{F_{21}}{F_{21}^{2} F_{22}} & \frac{F_{21}+F_{11}}{F_{11}^{2} F_{22}^{2}}
\end{array}\right]
\end{aligned}
$$

for the three unknowns $F_{11}, F_{21}$, and $F_{22}$ in terms of the known $\Omega_{11}, \Omega_{12}$, and $\Omega_{22}$. Note that the identifying restrictions are that $D=I$ (three restrictions) and $F_{12}=0$ (one restriction). (This system is just four nonlinear equations in three unknown - one of the equations for $\Omega_{12}$ is redundant. You do not need the Cholesky decomposition to solve it, since it could be solved with any numerical solver of non-linear equations-but why make life even more miserable?)

A practical consequence of this normalization is that the impulse response of shock $i$ equal to unity is exactly the same as the impulse response of shock $i$ equal to $\operatorname{Std}\left(u_{i t}\right)$ in the normalization in Section 11.7.2.

### 11.7.4 Other Identification Schemes*

## Reference: Bernanke (1986).

Not all economic models can be written in this recursive form. However, there are often cross-restrictions between different elements in $F$ or between elements in $F$ and $D$, or some other type of restrictions on $F$ which may allow us to identify the system.

Suppose we have (estimated) the parameters of the VAR (11.18), and that we want to impose $D=\operatorname{Cov}\left(u_{t}\right)=I$. From (11.19) we then have $(D=I)$

$$
\begin{equation*}
\Omega=F^{-1}\left(F^{-1}\right)^{\prime} \tag{11.22}
\end{equation*}
$$

As before we need $n(n-1) / 2$ restrictions on $F$, but this time we don't want to impose the restriction that all elements in $F$ above the principal diagonal are zero. Given these restrictions (whatever they are), we can solve for the remaining elements in $B$, typically with a numerical method for solving systems of non-linear equations.

### 11.7.5 What if the VAR Shocks are Uncorrelated $(\Omega=I)$ ?*

Suppose we estimate a VAR and find that the covariance matrix of the estimated residuals is (almost) an identity matrix (or diagonal). Does this mean that the identification is superfluous? No, not in general. Yes, if we also want to impose the restrictions that $F$ is triangular.

There are many ways to reshuffle the shocks and still get orthogonal shocks. Recall that the structural shocks are linear functions of the VAR shocks, $u_{t}=F \varepsilon_{t}$, and that we assume that $\operatorname{Cov}\left(\varepsilon_{t}\right)=\Omega=I$ and we want $\operatorname{Cov}\left(u_{t}\right)=I$, that, is from (11.19) we then have $(D=I)$

$$
\begin{equation*}
F F^{\prime}=I \tag{11.23}
\end{equation*}
$$

There are many such $F$ matrices: the class of those matrices even have a name: orthogonal matrices (all columns in $F$ are orthonormal). However, there is only one lower triangular $F$ which satisfies (11.23) (the one returned by a Cholesky decomposition, which is $I$ ).

Suppose you know that $F$ is lower triangular (and you intend to use this as the identifying assumption), but that your estimated $\Omega$ is (almost, at least) diagonal. The logic then requires that $F$ is not only lower triangular, but also diagonal. This means that $u_{t}=\varepsilon_{t}$ (up to a scaling factor). Therefore, a finding that the VAR shocks are uncorrelated combined with the identifying restriction that $F$ is triangular implies that the structural and reduced form shocks are proportional. We can draw no such conclusion if the identifying assumption is something else than lower triangularity

Example 25 (Rotation of vectors ("Givens rotations").) Consider the transformation of the vector $\varepsilon$ into the vector $u, u=G^{\prime} \varepsilon$, where $G=I_{n}$ except that $G_{i k}=c, G_{i k}=s$, $G_{k i}=-s$, and $G_{k k}=c$. If we let $c=\cos \theta$ and $s=\sin \theta$ for some angle $\theta$, then $G^{\prime} G=I$. To see this, consider the simple example where $i=2$ and $k=3$

$$
\left[\begin{array}{ccc}
1 & 0 & 0 \\
0 & c & s \\
0 & -s & c
\end{array}\right]^{\prime}\left[\begin{array}{ccc}
1 & 0 & 0 \\
0 & c & s \\
0 & -s & c
\end{array}\right]=\left[\begin{array}{ccc}
1 & 0 & 0 \\
0 & c^{2}+s^{2} & 0 \\
0 & 0 & c^{2}+s^{2}
\end{array}\right],
$$

which is an identity matrix since $\cos ^{2} \theta+\sin ^{2} \theta=1$. The transformation $u=G^{\prime} \varepsilon$ gives

$$
\begin{aligned}
u_{t} & =\varepsilon_{t} \text { for } t \neq i, k \\
u_{i} & =\varepsilon_{i} c-\varepsilon_{k} s \\
u_{k} & =\varepsilon_{i} s+\varepsilon_{k} c
\end{aligned}
$$

The effect of this transformation is to rotate the $i^{\text {th }}$ and $k^{\text {th }}$ vectors counterclockwise through an angle of $\theta$. (Try it in two dimensions.) There is an infinite number of such transformations (apply a sequence of such transformations with different $i$ and $k$, change $\theta$, etc.).

Example 26 (Givens rotations and the F matrix.) We could take $F$ in (11.23) to be (the transpose) of any such sequence of givens rotations. For instance, if $G_{1}$ and $G_{2}$ are givens rotations, then $F=G_{1}^{\prime}$ or $F=G_{2}^{\prime}$ or $F=G_{1}^{\prime} G_{2}^{\prime}$ are all valid.

### 11.7.6 Identification via Long-Run Restrictions - No Cointegration*

Suppose we have estimated a VAR system (11.1) for the first differences of some variables $y_{t}=\Delta x_{t}$, and that we have calculated the impulse response function as in (11.8), which we rewrite as

$$
\begin{align*}
\Delta x_{t} & =\varepsilon_{t}+C_{1} \varepsilon_{t-1}+C_{2} \varepsilon_{t-2}+\ldots \\
& =C(\mathrm{~L}) \varepsilon_{t}, \text { with } \operatorname{Cov}\left(\varepsilon_{t}\right)=\Omega \tag{11.24}
\end{align*}
$$

To find the MA of the level of $x_{t}$, we solve recursively

$$
\begin{align*}
& x_{t}=C(\mathrm{~L}) \varepsilon_{t}+x_{t-1} \\
&=C(\mathrm{~L}) \varepsilon_{t}+C(\mathrm{~L}) \varepsilon_{t-1}+x_{t-2} \\
& \vdots \\
&=C(\mathrm{~L})\left(\varepsilon_{t}+\varepsilon_{t-1}+\varepsilon_{t-2}+\ldots\right) \\
&=\varepsilon_{t}+\left(C_{1}+I\right) \varepsilon_{t-1}+\left(C_{2}+C_{1}+I\right) \varepsilon_{t-2}+\ldots  \tag{11.25}\\
&=C^{+}(\mathrm{L}) \varepsilon_{t}, \text { where } C_{s}^{+}=\sum_{j=0}^{s} C_{s} \text { with } C_{0}=I .
\end{align*}
$$

As before the structural shocks, $u_{t}$, are

$$
u_{t}=F \varepsilon_{t} \text { with } \operatorname{Cov}\left(u_{t}\right)=D
$$

The VMA in term of the structural shocks is therefore

$$
\begin{equation*}
x_{t}=C^{+}(\mathrm{L}) F^{-1} u_{t}, \text { where } C_{s}^{+}=\sum_{j=0}^{s} C_{s} \text { with } C_{0}=I \tag{11.26}
\end{equation*}
$$

The $C^{+}(\mathrm{L})$ polynomial is known from the estimation, so we need to identify $F$ in order to use this equation for impulse response function and variance decompositions with respect to the structural shocks.

As before we assume that $D=I$, so

$$
\begin{equation*}
\Omega=F^{-1} D\left(F^{-1}\right)^{\prime} \tag{11.27}
\end{equation*}
$$

in (11.19) gives $n(n+1) / 2$ restrictions.
We now add restrictions on the long run impulse responses. From (11.26) we have

$$
\begin{align*}
\lim _{s \rightarrow \infty} \frac{\partial x_{t+s}}{\partial u_{t}^{\prime}} & =\lim _{s \rightarrow \infty} C_{s}^{+} F^{-1} \\
& =C(1) F^{-1} \tag{11.28}
\end{align*}
$$

where $C(1)=\sum_{j=0}^{\infty} C_{s}$. We impose $n(n-1) / 2$ restrictions on these long run responses. Together we have $n^{2}$ restrictions, which allows to identify all elements in $F$.

In general, (11.27) and (11.28) is a set of non-linear equations which have to solved for the elements in $F$. However, it is common to assume that (11.28) is a lower triangular matrix. We can then use the following "trick" to find $F$. Since $\varepsilon_{t}=F^{-1} u_{t}$

$$
\begin{align*}
\mathrm{E} C(1) \varepsilon_{t} \varepsilon_{t}^{\prime} C(1)^{\prime} & =\mathrm{E} C(1) F^{-1} u_{t} u_{t}^{\prime}\left(F^{-1}\right)^{\prime} C(1)^{\prime} \\
C(1) \Omega C(1)^{\prime} & =C(1) F^{-1}\left(F^{-1}\right)^{\prime} C(1)^{\prime} . \tag{11.29}
\end{align*}
$$

We can therefore solve for a lower triangular matrix

$$
\begin{equation*}
\Lambda=C(1) F^{-1} \tag{11.30}
\end{equation*}
$$

by calculating the Cholesky decomposition of the left hand side of (11.29) (which is
available from the VAR estimate). Finally, we solve for $F^{-1}$ from (11.30).
Example 27 (The $2 \times 1$ case.) Suppose the structural form is

$$
\left[\begin{array}{ll}
F_{11} & F_{12} \\
F_{21} & F_{22}
\end{array}\right]\left[\begin{array}{l}
\Delta x_{t} \\
\Delta z_{t}
\end{array}\right]=\left[\begin{array}{ll}
B_{11} & B_{12} \\
B_{21} & B_{22}
\end{array}\right]\left[\begin{array}{l}
\Delta x_{t-1} \\
\Delta z_{t-1}
\end{array}\right]+\left[\begin{array}{l}
u_{1, t} \\
u_{2, t}
\end{array}\right] .
$$

and we have an estimate of the reduced form

$$
\left[\begin{array}{c}
\Delta x_{t} \\
\Delta z_{t}
\end{array}\right]=A\left[\begin{array}{l}
\Delta x_{t-1} \\
\Delta z_{t-1}
\end{array}\right]+\left[\begin{array}{l}
\varepsilon_{1, t} \\
\varepsilon_{2, t}
\end{array}\right] \text {, with } \operatorname{Cov}\left(\left[\begin{array}{l}
\varepsilon_{1, t} \\
\varepsilon_{2, t}
\end{array}\right]\right)=\Omega
$$

The VMA form (as in (11.24))

$$
\left[\begin{array}{l}
\Delta x_{t} \\
\Delta z_{t}
\end{array}\right]=\left[\begin{array}{l}
\varepsilon_{1, t} \\
\varepsilon_{2, t}
\end{array}\right]+A\left[\begin{array}{l}
\varepsilon_{1, t-1} \\
\varepsilon_{2, t-1}
\end{array}\right]+A^{2}\left[\begin{array}{l}
\varepsilon_{1, t-2} \\
\varepsilon_{2, t-2}
\end{array}\right]+\ldots
$$

and for the level (as in (11.25))

$$
\left[\begin{array}{l}
x_{t} \\
z_{t}
\end{array}\right]=\left[\begin{array}{l}
\varepsilon_{1, t} \\
\varepsilon_{2, t}
\end{array}\right]+(A+I)\left[\begin{array}{l}
\varepsilon_{1, t-1} \\
\varepsilon_{2, t-1}
\end{array}\right]+\left(A^{2}+A+I\right)\left[\begin{array}{l}
\varepsilon_{1, t-2} \\
\varepsilon_{2, t-2}
\end{array}\right]+\ldots
$$

or since $\varepsilon_{t}=F^{-1} u_{t}$
$\left[\begin{array}{l}x_{t} \\ z_{t}\end{array}\right]=F^{-1}\left[\begin{array}{l}u_{1, t} \\ u_{2, t}\end{array}\right]+(A+I) F^{-1}\left[\begin{array}{l}u_{1, t-1} \\ u_{2, t-1}\end{array}\right]+\left(A^{2}+A+I\right) F^{-1}\left[\begin{array}{l}u_{1, t-2} \\ u_{2, t-2}\end{array}\right]+\ldots$
There are $8+3$ parameters in the structural form and $4+3$ parameters in the VAR, so we need four restrictions. Assume that $\operatorname{Cov}\left(u_{t}\right)=I$ (three restrictions) and that the long run response of $u_{1, t-s}$ on $x_{t}$ is zero, that is,

$$
\begin{aligned}
{\left[\begin{array}{ll}
\text { unrestricted } & 0 \\
\text { unrestricted unrestricted }
\end{array}\right] } & =\left(I+A+A^{2}+\ldots\right)\left[\begin{array}{ll}
F_{11} & F_{12} \\
F_{21} & F_{22}
\end{array}\right]^{-1} \\
& =(I-A)^{-1}\left[\begin{array}{ll}
F_{11} & F_{12} \\
F_{21} & F_{22}
\end{array}\right]^{-1} \\
& =\left[\begin{array}{ll}
1-A_{11} & -A_{12} \\
-A_{21} & 1-A_{22}
\end{array}\right]^{-1}\left[\begin{array}{ll}
F_{11} & F_{12} \\
F_{21} & F_{22}
\end{array}\right]^{-1}
\end{aligned}
$$

The upper right element of the right hand side is

$$
\frac{-F_{12}+F_{12} A_{22}+A_{12} F_{11}}{\left(1-A_{22}-A_{11}+A_{11} A_{22}-A_{12} A_{21}\right)\left(F_{11} F_{22}-F_{12} F_{21}\right)}
$$

which is one restriction on the elements in $F$. The other three are given by $F^{-1}\left(F^{-1}\right)^{\prime}=$ $\Omega$, that is,

$$
\left[\begin{array}{cc}
\frac{F_{22}^{2}+F_{12}^{2}}{\left(F_{11} F_{22}-F_{12} F_{21}\right)^{2}} & -\frac{F_{22} F_{21}+F_{12} F_{11}}{\left(F_{11} F_{22}-F_{12} F_{21}\right)^{2}} \\
-\frac{F_{22} F_{21}+F_{12} F_{11}}{\left(F_{11} F_{22}-F_{12} F_{21}\right)^{2}} & \frac{F_{21}^{2}+F_{11}}{\left(F_{11} F_{22}-F_{12} F_{21}\right)^{2}}
\end{array}\right]=\left[\begin{array}{ll}
\Omega_{11} & \Omega_{12} \\
\Omega_{12} & \Omega_{22}
\end{array}\right] .
$$

### 11.8 Cointegration, Common Trends, and Identification via LongRun Restrictions*

These notes are a reading guide to Mellander, Vredin, and Warne (1992), which is well beyond the first year course in econometrics. See also Englund, Vredin, and Warne (1994). (I have not yet double checked this section.)

### 11.8.1 Common Trends Representation and Cointegration

The common trends representation of the $n$ variables in $y_{t}$ is

$$
\begin{align*}
y_{t} & =y_{0}+\Upsilon \tau_{t}+\Phi(\mathrm{L})\left[\begin{array}{l}
\varphi_{t} \\
\psi_{t}
\end{array}\right], \text { with } \operatorname{Cov}\left(\left[\begin{array}{l}
\varphi_{t} \\
\psi_{t}
\end{array}\right]\right)=I_{n}  \tag{11.31}\\
\tau_{t} & =\tau_{t-1}+\varphi_{t} \tag{11.32}
\end{align*}
$$

where $\Phi(\mathrm{L})$ is a stable matrix polynomial in the lag operator. We see that the $k \times 1$ vector $\varphi_{t}$ has permanent effects on (at least some elements in) $y_{t}$, while the $r \times 1(r=n-k) \psi_{t}$ does not.

The last component in (11.31) is stationary, but $\tau_{t}$ is a $k \times 1$ vector of random walks, so the $n \times k$ matrix $\Upsilon$ makes $y_{t}$ share the non-stationary components: there are $k$ common trends. If $k<n$, then we could find (at least) $r$ linear combinations of $y_{t}, \alpha^{\prime} y_{t}$ where $\alpha^{\prime}$ is an $r \times n$ matrix of cointegrating vectors, which are such that the trends cancel each other $\left(\alpha^{\prime} \Upsilon=\mathbf{0}\right)$.

Remark 28 (Lag operator.) We have the following rules: (i) $L^{k} x_{t}=x_{t-k}$; (ii) if $\Phi(L)=$
$a+b L^{-m}+c L^{n}$, then $\Phi(L)\left(x_{t}+y_{t}\right)=a\left(x_{t}+y_{t}\right)+b\left(x_{t+m}+y_{t+m}\right)+c\left(x_{t-n}+y_{t-n}\right)$ and $\Phi(1)=a+b+c$.

Example 29 (Söderlind and Vredin (1996)). Suppose we have
$y_{t}=\left[\begin{array}{l}\ln Y_{t} \text { (output) } \\ \ln P_{t} \text { (price level) } \\ \ln M_{t} \text { (money stock) } \\ \ln R_{t} \text { (gross interest rate) }\end{array}\right], \Upsilon=\left[\begin{array}{ll}0 & 1 \\ 1 & -1 \\ 1 & 0 \\ 0 & 0\end{array}\right]$, and $\tau_{t}=\left[\begin{array}{l}\text { money supply trend } \\ \text { productivity trend }\end{array}\right]$,
then we see that $\ln R_{t}$ and $\ln Y_{t}+\ln P_{t}-\ln M_{t}$ (that is, log velocity) are stationary, so

$$
\alpha^{\prime}=\left[\begin{array}{llll}
0 & 0 & 0 & 1 \\
1 & 1 & -1 & 0
\end{array}\right]
$$

are (or rather, span the space of) cointegrating vectors. We also see that $\alpha^{\prime} \Upsilon=\mathbf{0}_{2 \times 2}$.

### 11.8.2 VAR Representation

The VAR representation is as in (11.1). In practice, we often estimate the parameters in $A_{s}^{*}, \alpha$, the $n \times r$ matrix $\gamma$, and $\Omega=\operatorname{Cov}\left(\varepsilon_{t}\right)$ in the vector "error correction form"

$$
\begin{equation*}
\Delta y_{t}=A_{1}^{*} \Delta y_{t}+\ldots+A_{p-1}^{*} \Delta y_{t-p+1}+\gamma \alpha^{\prime} y_{t-1}+\varepsilon_{t}, \text { with } \operatorname{Cov}\left(\varepsilon_{t}\right)=\Omega \tag{11.33}
\end{equation*}
$$

This can easily be rewritten on the VAR form (11.1) or on the vector MA representation for $\Delta y_{t}$

$$
\begin{align*}
\Delta y_{t} & =\varepsilon_{t}+C_{1} \varepsilon_{t-1}+C_{2} \varepsilon_{t-2}+\ldots  \tag{11.34}\\
& =C(\mathrm{~L}) \varepsilon_{t} \tag{11.35}
\end{align*}
$$

To find the MA of the level of $y_{t}$, we recurse on (11.35)

$$
\begin{align*}
y_{t} & =C(\mathrm{~L}) \varepsilon_{t}+y_{t-1} \\
& =C(\mathrm{~L}) \varepsilon_{t}+C(\mathrm{~L}) \varepsilon_{t-1}+y_{t-2} \\
& \vdots  \tag{11.36}\\
& =C(\mathrm{~L})\left(\varepsilon_{t}+\varepsilon_{t-1}+\varepsilon_{t-2}+\ldots+\varepsilon_{0}\right)+y_{0} .
\end{align*}
$$

We now try to write (11.36) in a form which resembles the common trends representation
(11.31)-(11.32) as much as possible.

### 11.8.3 Multivariate Beveridge-Nelson decomposition

We want to split a vector of non-stationary series into some random walks and the rest (which is stationary). Rewrite (11.36) by adding and subtracting $C(1)\left(\varepsilon_{t}+\varepsilon_{t-1}+\ldots\right)$

$$
\begin{equation*}
y_{t}=C(1)\left(\varepsilon_{t}+\varepsilon_{t-1}+\varepsilon_{t-2}+\ldots+\varepsilon_{0}\right)+[C(\mathrm{~L})-C(1)]\left(\varepsilon_{t}+\varepsilon_{t-1}+\varepsilon_{t-2}+\ldots+\varepsilon_{0}\right) . \tag{11.37}
\end{equation*}
$$

Suppose $\varepsilon_{s}=0$ for $s<0$ and consider the second term in (11.37). It can be written

$$
\begin{align*}
& {\left[I+C_{1} \mathrm{~L}+C_{2} \mathrm{~L}^{2}+\ldots .-C(1)\right]\left(\varepsilon_{t}+\varepsilon_{t-1}+\varepsilon_{t-2}+\ldots+\varepsilon_{0}\right)} \\
& =/ * \text { since } C(1)=I+C_{1}+C_{2}+\ldots * / \\
& {\left[-C_{1}-C_{2}-C_{3}-\ldots\right] \varepsilon_{t}+\left[-C_{2}-C_{3}-\ldots\right] \varepsilon_{t-1}+\left[-C_{3}-\ldots\right] \varepsilon_{t-2}} \tag{11.38}
\end{align*}
$$

Now define the random walks

$$
\begin{align*}
\xi_{t} & =\xi_{t-1}+\varepsilon_{t}  \tag{11.39}\\
& =\varepsilon_{t}+\varepsilon_{t-1}+\varepsilon_{t-2}+\ldots+\varepsilon_{0} .
\end{align*}
$$

Use (11.38) and (11.39) to rewrite (11.37) as

$$
\begin{align*}
y_{t} & =C(1) \xi_{t}+C^{*}(\mathrm{~L}) \varepsilon_{t}, \text { where }  \tag{11.40}\\
C_{s}^{*} & =-\sum_{j=s+1}^{\infty} C_{j} \tag{11.41}
\end{align*}
$$

### 11.8.4 Identification of the Common Trends Shocks

Rewrite (11.31)-(11.32) and (11.39)-(11.40) as

$$
\begin{align*}
y_{t} & =C(1) \sum_{s=0}^{t} \varepsilon_{t}+C^{*}(\mathrm{~L}) \varepsilon_{t}, \text { with } \operatorname{Cov}\left(\varepsilon_{t}\right)=\Omega, \text { and }  \tag{11.42}\\
& =\left[\begin{array}{ll}
\Upsilon & \mathbf{0}_{n \times r}
\end{array}\right]\left[\begin{array}{l}
\sum_{s=0}^{t} \varphi_{t} \\
\psi_{t}
\end{array}\right]+\Phi(\mathrm{L})\left[\begin{array}{c}
\varphi_{t} \\
\psi_{t}
\end{array}\right], \text { with } \operatorname{Cov}\left(\left[\begin{array}{c}
\varphi_{t} \\
\psi_{t}
\end{array}\right]\right)=I_{n} . \tag{11.43}
\end{align*}
$$

Since both $\varepsilon_{t}$ and $\left[\begin{array}{cc}\varphi_{t}^{\prime} & \psi_{t}^{\prime}\end{array}\right]^{\prime}$ are white noise, we notice that the response of $y_{t+s}$ to either must be the same, that is,

$$
\left(C(1)+C_{s}^{*}\right) \varepsilon_{t}=\left(\left[\begin{array}{cc}
\Upsilon & \mathbf{0}_{n \times r}
\end{array}\right]+\Phi_{s}\right)\left[\begin{array}{l}
\varphi_{t}  \tag{11.44}\\
\psi_{t}
\end{array}\right] \text { for all } t \text { and } s \geq 0 .
$$

This means that the VAR shocks are linear combinations of the structural shocks (as in the standard setup without cointegration)

$$
\begin{align*}
{\left[\begin{array}{c}
\varphi_{t} \\
\psi_{t}
\end{array}\right] } & =F \varepsilon_{t} \\
& =\left[\begin{array}{c}
F_{k} \\
F_{r}
\end{array}\right] \varepsilon_{t} \tag{11.45}
\end{align*}
$$

Combining (11.44) and (11.45) gives that

$$
C(1)+C_{s}^{*}=\Upsilon F_{k}+\Phi_{s}\left[\begin{array}{l}
F_{k}  \tag{11.46}\\
F_{r}
\end{array}\right]
$$

must hold for all $s \geq 0$. In particular, it must hold for $s \rightarrow \infty$ where both $C_{s}^{*}$ and $\Phi_{s}$ vanishes

$$
\begin{equation*}
C(1)=\Upsilon F_{k} \tag{11.47}
\end{equation*}
$$

The identification therefore amounts to finding the $n^{2}$ coefficients in $F$, exactly as in the usual case without cointegration. Once that is done, we can calculate the impulse responses and variance decompositions with respect to the structural shocks by using
$\varepsilon_{t}=F^{-1}\left[\begin{array}{cc}\varphi_{t}^{\prime} & \psi_{t}^{\prime}\end{array}\right]^{\prime}$ in (11.42). ${ }^{2}$ As before, assumptions about the covariance matrix of the structural shocks are not enough to achieve identification. In this case, we typically rely on the information about long-run behavior (as opposed to short-run correlations) to supply the remaining restrictions.

- Step 1. From (11.31) we see that $\alpha^{\prime} \Upsilon=\mathbf{0}_{r \times k}$ must hold for $\alpha^{\prime} y_{t}$ to be stationary. Given an (estimate of) $\alpha$, this gives $r k$ equations from which we can identify $r k$ elements in $\Upsilon$. (It will soon be clear why it is useful to know $\Upsilon$ ).
- Step 2. From (11.44) we have $\Upsilon \varphi_{t}=C$ (1) $\varepsilon_{t}$ as $s \rightarrow \infty$. The variances of both sides must be equal

$$
\begin{array}{r}
\mathrm{E} \Upsilon \varphi_{t} \varphi_{t}^{\prime} \Upsilon^{\prime}=\mathrm{E} C(1) \varepsilon_{t} \varepsilon_{t}^{\prime} C(1)^{\prime}, \text { or } \\
\Upsilon \Upsilon^{\prime}=C(1) \Omega C(1)^{\prime}, \tag{11.48}
\end{array}
$$

which gives $k(k+1) / 2$ restrictions on $\Upsilon$ (the number of unique elements in the symmetric $\Upsilon \Upsilon^{\prime}$ ). (However, each column of $\Upsilon$ is only identified up to a sign transformation: neither step 1 or 2 is affected by multiplying each element in column $j$ of $\Upsilon$ by -1.)

- Step 3. $\Upsilon$ has $n k$ elements, so we still need $n k-r k-k(k+1) / 2=k(k-1) / 2$ further restrictions on $\Upsilon$ to identify all elements. They could be, for instance, that money supply shocks have no long run effect on output (some $\Upsilon_{i j}=0$ ). We now know $\Upsilon$.
- Step 4. Combining $\operatorname{Cov}\left(\left[\begin{array}{l}\varphi_{t} \\ \psi_{t}\end{array}\right]\right)=I_{n}$ with (11.45) gives

$$
\left[\begin{array}{cc}
I_{k} & 0  \tag{11.49}\\
0 & I_{r}
\end{array}\right]=\left[\begin{array}{c}
F_{k} \\
F_{r}
\end{array}\right] \Omega\left[\begin{array}{c}
F_{k} \\
F_{r}
\end{array}\right]^{\prime}
$$

which gives $n(n+1) / 2$ restrictions.

- Step 4a. Premultiply (11.47) with $\Upsilon^{\prime}$ and solve for $F_{k}$

$$
\begin{equation*}
F_{k}=\left(\Upsilon^{\prime} \Upsilon\right)^{-1} \Upsilon^{\prime} C(1) \tag{11.50}
\end{equation*}
$$

${ }^{2}$ Equivalently, we can use (11.47) and (11.46) to calculate $\Upsilon$ and $\Phi_{s}$ (for all $s$ ) and then calculate the impulse response function from (11.43)
(This means that $\mathrm{E} \varphi_{t} \varphi_{t}^{\prime}=F_{k} \Omega F_{k}^{\prime}=\left(\Upsilon^{\prime} \Upsilon\right)^{-1} \Upsilon^{\prime} C(1) \Omega C(1)^{\prime} \Upsilon\left(\Upsilon^{\prime} \Upsilon\right)^{-1}$. From (11.48) we see that this indeed is $I_{k}$ as required by (11.49).) We still need to identify $F_{r}$.

- Step 4b. From (11.49), $\mathrm{E} \varphi_{t} \psi_{t}^{\prime}=\mathbf{0}_{k \times r}$, we get $F_{k} \Omega F_{r}^{\prime}=\mathbf{0}_{k \times r}$, which gives $k r$ restrictions on the $r n$ elements in $F_{r}$. Similarly, from $\mathrm{E} \psi_{t} \psi_{t}^{\prime}=I_{r}$, we get $F_{r} \Omega F_{r}^{\prime}=I_{r}$, which gives $r(r+1) / 2$ additional restrictions on $F_{r}$. We still need $r(r-1) / 2$ restrictions. Exactly how they look does not matter for the impulse response function of $\varphi_{t}$ (as long as $\mathrm{E} \varphi_{t} \psi_{t}^{\prime}=\mathbf{0}$ ). Note that restrictions on $F_{r}$ are restrictions on $\partial y_{t} / \partial \psi_{t}^{\prime}$, that is, on the contemporaneous response. This is exactly as in the standard case without cointegration.

A summary of identifying assumptions used by different authors is found in Englund, Vredin, and Warne (1994).

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## 12 Kalman filter

### 12.1 Conditional Expectations in a Multivariate Normal Distribution

Reference: Harvey (1989), Lütkepohl (1993), and Hamilton (1994)
Suppose $Z_{m \times 1}$ and $X_{n \times 1}$ are jointly normally distributed

$$
\left[\begin{array}{c}
Z  \tag{12.1}\\
X
\end{array}\right]=N\left(\left[\begin{array}{c}
\bar{Z} \\
\bar{X}
\end{array}\right],\left[\begin{array}{cc}
\Sigma_{z z} & \Sigma_{z x} \\
\Sigma_{x z} & \Sigma_{x x}
\end{array}\right]\right) .
$$

The distribution of the random variable $Z$ conditional on that $X=x$ is also normal with mean (expectation of the random variable $Z$ conditional on that the random variable $X$ has the value $x$ )

$$
\begin{equation*}
\underbrace{\mathrm{E}(Z \mid x)}_{m \times 1}=\underbrace{\bar{Z}}_{m \times 1}+\underbrace{\Sigma_{z x} x}_{m \times n} \underbrace{\Sigma_{x x}^{-1}}_{n \times n} \underbrace{(x-\bar{X})}_{n \times 1}, \tag{12.2}
\end{equation*}
$$

and variance (variance of $Z$ conditional on that $X=x$ )

$$
\begin{align*}
\operatorname{Var}(Z \mid x) & =\mathrm{E}\left\{[Z-\mathrm{E}(Z \mid x)]^{2} \mid x\right\} \\
& =\Sigma_{z z}-\Sigma_{z x} \Sigma_{x x}^{-1} \Sigma_{x z} . \tag{12.3}
\end{align*}
$$

The conditional variance is the variance of the prediction error $Z-\mathrm{E}(Z \mid x)$.
Both $\mathrm{E}(Z \mid x)$ and $\operatorname{Var}(Z \mid x)$ are in general stochastic variables, but for the multivariate normal distribution $\operatorname{Var}(Z \mid x)$ is constant. Note that $\operatorname{Var}(Z \mid x)$ is less than $\Sigma_{z z}$ (in a matrix sense) if $x$ contains any relevant information (so $\Sigma_{z x}$ is not zero, that is, $\mathrm{E}(z \mid x)$ is not a constant).

It can also be useful to know that $\operatorname{Var}(Z)=\mathrm{E}[\operatorname{Var}(Z \mid X)]+\operatorname{Var}[\mathrm{E}(Z \mid X)]$ (the $X$ is now random), which here becomes $\Sigma_{z z}-\Sigma_{z x} \Sigma_{x x}^{-1} \Sigma_{x z}+\Sigma_{z x} \Sigma_{x x}^{-1} \operatorname{Var}(X) \Sigma_{x x}^{-1} \Sigma_{x z}=\Sigma_{z z}$.

### 12.2 Kalman Recursions

### 12.2.1 State space form

The measurement equation is

$$
\begin{equation*}
y_{t}=Z \alpha_{t}+\epsilon_{t}, \text { with } \operatorname{Var}\left(\epsilon_{t}\right)=H \tag{12.4}
\end{equation*}
$$

where $y_{t}$ and $\epsilon_{t}$ are $n \times 1$ vectors, and $Z$ an $n \times m$ matrix. (12.4) expresses some observable variables $y_{t}$ in terms of some (partly) unobservable state variables $\alpha_{t}$. The transition equation for the states is

$$
\begin{equation*}
\alpha_{t}=T \alpha_{t-1}+u_{t}, \text { with } \operatorname{Var}\left(u_{t}\right)=Q, \tag{12.5}
\end{equation*}
$$

where $\alpha_{t}$ and $u_{t}$ are $m \times 1$ vectors, and $T$ an $m \times m$ matrix. This system is time invariant since all coefficients are constant. It is assumed that all errors are normally distributed, and that $\mathrm{E}\left(\epsilon_{t} u_{t-s}\right)=0$ for all $s$.

Example $1(A R(2)$.$) The process x_{t}=\rho_{1} x_{t-1}+\rho_{2} x_{t-2}+e_{t}$ can be rewritten as

$$
\begin{gathered}
\underbrace{x_{t}}_{y_{t}}=\underbrace{\left[\begin{array}{ll}
1 & 0
\end{array}\right]}_{Z} \underbrace{\left[\begin{array}{c}
x_{t} \\
x_{t-1}
\end{array}\right]}_{\alpha_{t}}+\underbrace{0}_{\epsilon_{t}}, \\
\underbrace{\left[\begin{array}{c}
x_{t} \\
x_{t-1}
\end{array}\right]}_{\alpha_{t}}=\underbrace{\left[\begin{array}{cc}
\rho_{1} & \rho_{2} \\
1 & 0
\end{array}\right]}_{T} \underbrace{\left[\begin{array}{c}
x_{t-1} \\
x_{t-2}
\end{array}\right]}_{\alpha_{t-1}}+\underbrace{\left[\begin{array}{c}
e_{t} \\
0
\end{array}\right]}_{u_{t}},
\end{gathered}
$$

with $H=0$, and $Q=\left[\begin{array}{cc}\operatorname{Var}\left(e_{t}\right) & 0 \\ 0 & 0\end{array}\right]$. In this case $n=1, m=2$.

### 12.2.2 Prediction equations: $\mathbf{E}\left(\alpha_{t} \mid \mathbf{I}_{t-1}\right)$

Suppose we have an estimate of the state in $t-1$ based on the information set in $t-1$, denoted by $\hat{\alpha}_{t-1}$, and that this estimate has the variance

$$
\begin{equation*}
P_{t-1}=\mathrm{E}\left[\left(\hat{\alpha}_{t-1}-\alpha_{t-1}\right)\left(\hat{\alpha}_{t-1}-\alpha_{t-1}\right)^{\prime}\right] . \tag{12.6}
\end{equation*}
$$

Now we want an estimate of $\alpha_{t}$ based $\hat{\alpha}_{t-1}$. From (12.5) the obvious estimate, denoted by $\alpha_{t \mid t-1}$, is

$$
\begin{equation*}
\hat{\alpha}_{t \mid t-1}=T \hat{\alpha}_{t-1} \tag{12.7}
\end{equation*}
$$

The variance of the prediction error is

$$
\begin{align*}
P_{t \mid t-1} & =\mathrm{E}\left[\left(\alpha_{t}-\hat{\alpha}_{t \mid t-1}\right)\left(\alpha_{t}-\hat{\alpha}_{t \mid t-1}\right)^{\prime}\right] \\
& =\mathrm{E}\left\{\left[T \alpha_{t-1}+u_{t}-T \hat{\alpha}_{t-1}\right]\left[T \alpha_{t-1}+u_{t}-T \hat{\alpha}_{t-1}\right]^{\prime}\right\} \\
& =\mathrm{E}\left\{\left[T\left(\hat{\alpha}_{t-1}-\alpha_{t-1}\right)-u_{t}\right]\left[T\left(\hat{\alpha}_{t-1}-\alpha_{t-1}\right)-u_{t}\right]^{\prime}\right\} \\
& =T \mathrm{E}\left[\left(\hat{\alpha}_{t-1}-\alpha_{t-1}\right)\left(\hat{\alpha}_{t-1}-\alpha_{t-1}\right)^{\prime}\right] T^{\prime}+\mathrm{E} u_{t} u_{t}^{\prime} \\
& =T P_{t-1} T^{\prime}+Q, \tag{12.8}
\end{align*}
$$

where we have used (12.5), (12.6), and the fact that $u_{t}$ is uncorrelated with $\hat{\alpha}_{t-1}-\alpha_{t-1}$.
Example 2 (AR(2) continued.) By substitution we get

$$
\begin{aligned}
& \hat{\alpha}_{t \mid t-1}=\left[\begin{array}{c}
\hat{x}_{t \mid t-1} \\
\hat{x}_{t-1 \mid t-1}
\end{array}\right]=\left[\begin{array}{cc}
\rho_{1} & \rho_{2} \\
1 & 0
\end{array}\right]\left[\begin{array}{l}
\hat{x}_{t-1 \mid t-1} \\
\hat{x}_{t-2 \mid t-1}
\end{array}\right], \text { and } \\
& P_{t \mid t-1}=\left[\begin{array}{cc}
\rho_{1} & \rho_{2} \\
1 & 0
\end{array}\right] P_{t-1}\left[\begin{array}{ll}
\rho_{1} & 1 \\
\rho_{2} & 0
\end{array}\right]+\left[\begin{array}{cc}
\operatorname{Var}\left(\epsilon_{t}\right) & 0 \\
0 & 0
\end{array}\right]
\end{aligned}
$$

If we treat $x_{-1}$ and $x_{0}$ as given, then $P_{0}=\mathbf{0}_{2 \times 2}$ which would give $P_{1 \mid 0}=\left[\begin{array}{cc}\operatorname{Var}\left(\epsilon_{t}\right) & 0 \\ 0 & 0\end{array}\right]$
12.2.3 Updating equations: $\mathbf{E}\left(\alpha_{t} \mid \mathbf{I}_{t-1}\right) \rightarrow \mathbf{E}\left(\alpha_{t} \mid \mathbf{I}_{t}\right)$

The best estimate of $y_{t}$, given $\hat{a}_{t \mid t-1}$, follows directly from (12.4)

$$
\begin{equation*}
\hat{y}_{t \mid t-1}=Z \hat{\alpha}_{t \mid t-1}, \tag{12.9}
\end{equation*}
$$

with prediction error

$$
\begin{equation*}
v_{t}=y_{t}-\hat{y}_{t \mid t-1}=Z\left(\alpha_{t}-\hat{\alpha}_{t \mid t-1}\right)+\epsilon_{t} \tag{12.10}
\end{equation*}
$$

The variance of the prediction error is

$$
\begin{align*}
F_{t} & =\mathrm{E}\left(v_{t} v_{t}^{\prime}\right) \\
& =\mathrm{E}\left\{\left[Z\left(\alpha_{t}-\hat{\alpha}_{t \mid t-1}\right)+\epsilon_{t}\right]\left[Z\left(\alpha_{t}-\hat{\alpha}_{t \mid t-1}\right)+\epsilon_{t}\right]^{\prime}\right\} \\
& =Z \mathrm{E}\left[\left(\alpha_{t}-\hat{\alpha}_{t \mid t-1}\right)\left(\alpha_{t}-\hat{\alpha}_{t \mid t-1}\right)^{\prime}\right] Z^{\prime}+\mathrm{E} \epsilon_{t} \epsilon_{t}^{\prime} \\
& =Z P_{t \mid t-1} Z^{\prime}+H, \tag{12.11}
\end{align*}
$$

where we have used the definition of $P_{t \mid t-1}$ in (12.8), and of $H$ in 12.4. Similarly, the covariance of the prediction errors for $y_{t}$ and for $\alpha_{t}$ is

$$
\begin{align*}
\operatorname{Cov}\left(\alpha_{t}-\hat{\alpha}_{t \mid t-1}, y_{t}-\hat{y}_{t \mid t-1}\right) & =\mathrm{E}\left(\alpha_{t}-\hat{\alpha}_{t \mid t-1}\right)\left(y_{t}-\hat{y}_{t \mid t-1}\right) \\
& =\mathrm{E}\left\{\left(\alpha_{t}-\hat{\alpha}_{t \mid t-1}\right)\left[Z\left(\alpha_{t}-\hat{\alpha}_{t \mid t-1}\right)+\epsilon_{t}\right]^{\prime}\right\} \\
& =\mathrm{E}\left[\left(\alpha_{t}-\hat{\alpha}_{t \mid t-1}\right)\left(\alpha_{t}-\hat{\alpha}_{t \mid t-1}\right)^{\prime}\right] Z^{\prime} \\
& =P_{t \mid t-1} Z^{\prime} . \tag{12.12}
\end{align*}
$$

Suppose that $y_{t}$ is observed and that we want to update our estimate of $\alpha_{t}$ from $\hat{\alpha}_{t \mid t-1}$ to $\hat{\alpha}_{t}$, where we want to incorporate the new information conveyed by $y_{t}$.

Example 3 (AR(2) continued.) We get

$$
\begin{aligned}
& \hat{y}_{t \mid t-1}=Z \hat{\alpha}_{t \mid t-1}=\left[\begin{array}{ll}
1 & 0
\end{array}\right]\left[\begin{array}{c}
\hat{x}_{t \mid t-1} \\
\hat{x}_{t-1 \mid t-1}
\end{array}\right]=\hat{x}_{t \mid t-1}=\rho_{1} \hat{x}_{t-1 \mid t-1}+\rho_{2} \hat{x}_{t-2 \mid t-1}, \text { and } \\
& \qquad F_{t}=\left[\begin{array}{ll}
1 & 0
\end{array}\right]\left\{\left[\begin{array}{cc}
\rho_{1} & \rho_{2} \\
1 & 0
\end{array}\right] P_{t-1}\left[\begin{array}{cc}
\rho_{1} & 1 \\
\rho_{2} & 0
\end{array}\right]+\left[\begin{array}{cc}
\operatorname{Var}\left(\epsilon_{t}\right) & 0 \\
0 & 0
\end{array}\right]\right\}\left[\begin{array}{ll}
1 & 0
\end{array}\right]^{\prime} . \\
& \text { If } P_{0}=\mathbf{0}_{2 \times 2} \text { as before, then } F_{1}=P_{1}=\left[\begin{array}{ccc}
\operatorname{Var}\left(\epsilon_{t}\right) & 0 \\
0 & 0
\end{array}\right] .
\end{aligned}
$$

By applying the rules (12.2) and (12.3) we note that the expectation of $\alpha_{t}$ (like $z$ in (12.2)) conditional on $y_{t}$ (like $x$ in (12.2)) is (note that $y_{t}$ is observed so we can use it to guess $\alpha_{t}$ )

$$
\begin{equation*}
\underbrace{\hat{\alpha}_{t}}_{\mathrm{E}(z \mid x)}=\underbrace{\hat{\alpha}_{t \mid t-1}}_{\mathrm{E} z}+\underbrace{P_{t \mid t-1} Z^{\prime}}_{\Sigma_{z x}}(\underbrace{Z P_{t \mid t-1} Z^{\prime}+H}_{\Sigma_{x x}=F_{t}})^{-1}(y_{t}-\underbrace{Z \hat{\alpha}_{t \mid t-1}}_{\mathrm{E} x}) \tag{12.13}
\end{equation*}
$$

with variance

$$
\begin{equation*}
\underbrace{P_{t}}_{\operatorname{Var}(z \mid x)}=\underbrace{P_{t \mid t-1}}_{\Sigma_{z z}}-\underbrace{P_{t \mid t-1}^{\prime} Z^{\prime}}_{\Sigma_{z x}} \underbrace{\left(Z P_{t \mid t-1} Z^{\prime}+H\right)^{-1}}_{\Sigma_{x x}^{-1}} \underbrace{Z P_{t \mid t-1}}_{\Sigma_{x z}} \tag{12.14}
\end{equation*}
$$

where $\hat{\alpha}_{t \mid t-1}$ ("Ez") is from (12.7), $P_{t \mid t-1} Z^{\prime}$ (" $\Sigma_{z x}$ ") from (12.12), $Z P_{t \mid t-1} Z^{\prime}+H$ (" $\Sigma_{x x}$ ") from (12.11), and $Z \hat{\alpha}_{t \mid t-1}$ ("E $x$ ") from (12.9).
(12.13) uses the new information in $y_{t}$, that is, the observed prediction error, in order to update the estimate of $\alpha_{t}$ from $\hat{\alpha}_{t \mid t-1}$ to $\hat{\alpha}_{t}$.

Proof. The last term in (12.14) follows from the expected value of the square of the last term in (12.13)
$P_{t \mid t-1} Z^{\prime}\left(Z P_{t \mid t-1} Z^{\prime}+H\right)^{-1} \mathrm{E}\left(y_{t}-Z \alpha_{t \mid t-1}\right)\left(y_{t}-Z \alpha_{t \mid t-1}\right)^{\prime}\left(Z P_{t \mid t-1} Z^{\prime}+H\right)^{-1} Z P_{t \mid t-1}$,
where we have exploited the symmetry of covariance matrices. Note that $y_{t}-Z \alpha_{t \mid t-1}=$ $y_{t}-\hat{y}_{t \mid t-1}$, so the middle term in the previous expression is

$$
\begin{equation*}
\mathrm{E}\left(y_{t}-Z \alpha_{t \mid t-1}\right)\left(y_{t}-Z \alpha_{t \mid t-1}\right)^{\prime}=Z P_{t \mid t-1} Z^{\prime}+H \tag{12.16}
\end{equation*}
$$

Using this gives the last term in (12.14).

### 12.2.4 The Kalman Algorithm

The Kalman algorithm calculates optimal predictions of $\alpha_{t}$ in a recursive way. You can also calculate the prediction errors $v_{t}$ in (12.10) as a by-prodct, which turns out to be useful in estimation.

1. Pick starting values for $P_{0}$ and $\alpha_{0}$. Let $t=1$.
2. Calculate (12.7), (12.8), (12.13), and (12.14) in that order. This gives values for $\hat{\alpha}_{t}$ and $P_{t}$. If you want $v_{t}$ for estimation purposes, calculate also (12.10) and (12.11). Increase $t$ with one step.
3. Iterate on 2 until $t=T$.

One choice of starting values that work in stationary models is to set $P_{0}$ to the unconditional covariance matrix of $\alpha_{t}$, and $\alpha_{0}$ to the unconditional mean. This is the matrix $P$
to which (12.8) converges: $P=T P T^{\prime}+Q$. (The easiest way to calculate this is simply to start with $P=I$ and iterate until convergence.)

In non-stationary model we could set

$$
\begin{equation*}
P_{0}=1000 * I_{m}, \text { and } \alpha_{0}=\mathbf{0}_{m \times 1}, \tag{12.17}
\end{equation*}
$$

in which case the first $m$ observations of $\hat{\alpha}_{t}$ and $v_{t}$ should be disregarded.

### 12.2.5 MLE based on the Kalman filter

For any (conditionally) Gaussian time series model for the observable $y_{t}$ the log likelihood for an observation is

$$
\begin{equation*}
\ln L_{t}=-\frac{n}{2} \ln (2 \pi)-\frac{1}{2} \ln \left|F_{t}\right|-\frac{1}{2} v_{t}^{\prime} F_{t}^{-1} v_{t} . \tag{12.18}
\end{equation*}
$$

In case the starting conditions are as in (12.17), the overall log likelihood function is

$$
\ln L=\left\{\begin{array}{l}
\sum_{t=1}^{T} \ln L_{t} \text { in stationary models }  \tag{12.19}\\
\sum_{t=m+1}^{T} \ln L_{t} \text { in non-stationary models. }
\end{array}\right.
$$

### 12.2.6 Inference and Diagnostics

We can, of course, use all the asymptotic MLE theory, like likelihood ratio tests etc. For diagnostoic tests, we will most often want to study the normalized residuals

$$
\tilde{v}_{i t}=v_{i t} / \sqrt{\text { element } i i \text { in } F_{t}}, i=1, \ldots, n
$$

since element $i i$ in $F_{t}$ is the standard deviation of the scalar residual $v_{i t}$. Typical tests are CUSUMQ tests for structural breaks, various tests for serial correlation, heteroskedastic-
ity, and normality.

## Bibliography

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## 13 Outliers and Robust Estimators

### 13.1 Influential Observations and Standardized Residuals

Reference: Greene (2000) 6.9; Rousseeuw and Leroy (1987)
Consider the linear model

$$
\begin{equation*}
y_{t}=x_{t}^{\prime} \beta_{0}+u_{t}, \tag{13.1}
\end{equation*}
$$

where $x_{t}$ is $k \times 1$. The LS estimator

$$
\hat{\beta}=\left(\sum_{t=1}^{T} x_{t} x_{t}^{\prime}\right)^{-1} \sum_{t=1}^{T} x_{t} y_{t},
$$

which is the solution to

$$
\min _{\beta} \sum_{t=1}^{T}\left(y_{t}-x_{t}^{\prime} \beta\right)^{2}
$$

The fitted values and residuals are

$$
\hat{y}_{t}=x_{t}^{\prime} \hat{\beta}, \text { and } \hat{u}_{t}=y_{t}-\hat{y}_{t} .
$$

Suppose we were to reestimate $\beta$ on the whole sample, except observation $s$. This would give us an estimate $\hat{\beta}^{(s)}$. The fitted values and residual are then

$$
\begin{equation*}
\hat{y}_{t}^{(s)}=x_{t}^{\prime} \hat{\beta}^{(s)}, \text { and } \hat{u}_{t}^{(s)}=y_{t}-\hat{y}_{t}^{(s)} \tag{13.5}
\end{equation*}
$$

A common way to study the sensitivity of the results with respect to excluding observations is to plot $\hat{\beta}^{(s)}-\hat{\beta}$, and $\hat{y}_{s}^{(s)}-\hat{y}_{s}$. Note that we here plot the fitted value of $y_{s}$ using the coefficients estimated by excluding observation $s$ from the sample. Extreme values prompt a closer look at data (errors in data?) and perhaps also a more robust estimation method than LS, which is very sensitive to outliers.

Another useful way to spot outliers is to study the standardized residuals, $\hat{u}_{s} / \hat{\sigma}$ and $\hat{u}_{s}^{(s)} / \hat{\sigma}^{(s)}$, where $\hat{\sigma}$ and $\hat{\sigma}^{(s)}$ are standard deviations estimated from the whole sample and excluding observation $s$, respectively. Values below - 2 or above 2 warrant attention (recall

## hat $\operatorname{Pr}(x>1.96) \approx 0.025$ in a $N(0,1)$ distribution)

Sometimes the residuals are instead standardized by taking into account the uncertainty of the estimated coefficients. Note that

$$
\begin{align*}
\hat{u}_{t}^{(s)} & =y_{t}-x_{t}^{\prime} \hat{\beta}^{(s)} \\
& =u_{t}+x_{t}^{\prime}\left(\beta-\hat{\beta}^{(s)}\right), \tag{13.6}
\end{align*}
$$

since $y_{t}=x_{t}^{\prime} \beta+u_{t}$. The variance of $\hat{u}_{t}$ is therefore the variance of the sum on the right hand side of this expression. When we use the variance of $u_{t}$ as we did above to standardize the residuals, then we disregard the variance of $\hat{\beta}^{(s)}$. In general, we have

$$
\begin{equation*}
\operatorname{Var}\left(\hat{u}_{t}^{(s)}\right)=\operatorname{Var}\left(u_{t}\right)+x_{t}^{\prime} \operatorname{Var}\left(\beta-\hat{\beta}^{(s)}\right) x_{t}+2 \operatorname{Cov}\left[u_{t}, x_{t}^{\prime}\left(\beta-\hat{\beta}^{(s)}\right)\right] . \tag{13.7}
\end{equation*}
$$

When $t=s$, which is the case we care about, the covariance term drops out since $\hat{\beta}^{(s)}$ cannot be correlated with $u_{s}$ since period $s$ is not used in the estimation (this statement assumes that shocks are not autocorrelated). The first term is then estimated as the usual variance of the residuals (recall that period $s$ is not used) and the second term is the estimated covariance matrix of the parameter vector (once again excluding period $s$ ) preand postmultiplied by $x_{s}$.

Example 1 (Errors are iid independent of the regressors.) In this case the variance of the parameter vector is estimated as $\hat{\sigma}^{2}\left(\Sigma x_{t} x_{t}^{\prime}\right)^{-1}$ (excluding period $s$ ), so we have

$$
\operatorname{Var}\left(\hat{u}_{t}^{(s)}\right)=\hat{\sigma}^{2}\left(1+x_{s}^{\prime}\left(\Sigma x_{t} x_{t}^{\prime}\right)^{-1} x_{s}\right)
$$

### 13.2 Recursive Residuals*

## Reference: Greene (2000) 7.8

Recursive residuals are a version of the technique discussed in Section 13.1. They are used when data is a time series. Suppose we have a sample $t=1, \ldots, T$,.and that $t=1, \ldots, s$ are used to estimate a first estimate, $\hat{\beta}^{[s]}$ (not to be confused with $\hat{\beta}^{(s)}$ used in Section 13.1). We then make a one-period ahead forecast and record the fitted value and the forecast error

$$
\begin{equation*}
\hat{y}_{s+1}^{[s]}=x_{s+1}^{\prime} \hat{\beta}^{[s]}, \text { and } \hat{u}_{s+1}^{[s]}=y_{s+1}-\hat{y}_{s+1}^{[s]} . \tag{13.8}
\end{equation*}
$$

Rescursive residuals from $\operatorname{AR}(1)$ with corr=0.85USUM statistics and $95 \%$ confidence band



Figure 13.1: This figure shows recursive residuals and CUSUM statistics, when data are simulated from $y_{t}=0.85 y_{t-1}+u_{t}$, with $\operatorname{Var}\left(u_{t}\right)=1$.

This is repeated for the rest of the sample by extending the sample used in the estimation by one period, making a one-period ahead forecast, and then repeating until we reach the end of the sample.

A first diagnosis can be made by examining the standardized residuals, $\hat{u}_{s+1}^{[s]} / \hat{\sigma}^{[s]}$, where $\hat{\sigma}^{[s]}$ can be estimated as in (13.7) with a zero covariance term, since $u_{s+1}$ is not correlated with data for earlier periods (used in calculating $\hat{\beta}^{[s]}$ ), provided errors are not autocorrelated. As before, standardized residuals outside $\pm 2$ indicates problems: outliers or structural breaks (if the residuals are persistently outside $\pm 2$ ).

The CUSUM test uses these standardized residuals to form a sequence of test statistics. A (persistent) jump in the statistics is a good indicator of a structural break. Suppose we use $r$ observations to form the first estimate of $\beta$, so we calculate $\hat{\beta}^{[s]}$ and $\hat{u}_{s+1}^{[s]} / \hat{\sigma}^{[s]}$ for $s=r, \ldots, T$. Define the cumulative sums of standardized residuals

$$
\begin{equation*}
W_{t}=\sum_{s=r}^{t} \hat{u}_{s+1}^{[s]} / \hat{\sigma}^{[s]}, t=r, \ldots, T . \tag{13.9}
\end{equation*}
$$

Under the null hypothesis that no structural breaks occurs, that is, that the true $\beta$ is the same for the whole sample, $W_{t}$ has a zero mean and a variance equal to the number of elements in the sum, $t-r+1$. This follows from the fact that the standardized residuals all have zero mean and unit variance and are uncorrelated with each other. Typically, $W_{t}$ is plotted along with a $95 \%$ confidence interval, which can be shown to be $\pm(a \sqrt{T-r}+2 a(t-r) / \sqrt{T-r})$ with $a=0.948$. The hypothesis of no structural break is rejected if the $W_{t}$ is outside this band for at least one observation. (The derivation of this confidence band is somewhat tricky, but it incorporates the fact that $W_{t}$ and $W_{t+1}$

OLS vs LAD


Figure 13.2: This figure shows an example of how LS and LAD can differ. In this case $y_{t}=0.75 x_{t}+u_{t}$, but only one of the errors has a non-zero value.
are very correlated.)

### 13.3 Robust Estimation

Reference: Greene (2000) 9.8.1; Rousseeuw and Leroy (1987); Donald and Maddala (1993); and Judge, Griffiths, Lütkepohl, and Lee (1985) 20.4.

The idea of robust estimation is to give less weight to extreme observations than in Least Squares. When the errors are normally distributed, then there should be very few extreme observations, so LS makes a lot of sense (and is indeed the MLE). When the errors have distributions with fatter tails (like the Laplace or two-tailed exponential distribution, $f(u)=\exp (-|u| / \sigma) / 2 \sigma)$, then LS is no longer optimal and can be fairly sensitive to outliers. The ideal way to proceed would be to apply MLE, but the true distribution is often unknown. Instead, one of the "robust estimators" discussed below is often used.

Let $\hat{u}_{t}=y_{t}-x_{t}^{\prime} \hat{\beta}$. Then, the least absolute deviations (LAD), least median squares

LMS), and least trimmed squares (LTS) estimators solve

$$
\begin{align*}
& \hat{\beta}_{L A D}=\arg \min _{\beta} \sum_{t=1}^{T}\left|\hat{u}_{t}\right|  \tag{13.10}\\
& \hat{\beta}_{L M S}=\arg \min _{\beta}\left[\operatorname{median}\left(\hat{u}_{t}^{2}\right)\right]  \tag{13.11}\\
& \hat{\beta}_{L T S}=\arg \min _{\beta} \sum_{i=1}^{h} \hat{u}_{i}^{2}, \hat{u}_{1}^{2} \leq \hat{u}_{2}^{2} \leq \ldots \text { and } h \leq T . \tag{13.12}
\end{align*}
$$

Note that the LTS estimator in (13.12) minimizes of the sum of the $h$ smallest squared residuals.

These estimators involve non-linearities, so they are more computationally intensive than LS. In some cases, however, a simple iteration may work.

Example 2 (Algorithm for LAD.) The LAD estimator can be written

$$
\hat{\beta}_{L A D}=\arg \min _{\beta} \sum_{t=1}^{T} w_{t} \hat{u}_{t}^{2}, w_{t}=1 /\left|\hat{u}_{t}\right|
$$

so it is a weighted least squares where both $y_{t}$ and $x_{t}$ are multiplied by $1 /\left|\hat{u}_{t}\right|$. It can be shown that iterating on $L S$ with the weights given by $1 /\left|\hat{u}_{t}\right|$, where the residuals are from the previous iteration, converges very quickly to the LAD estimator.

It can be noted that LAD is actually the MLE for the Laplace distribution discussed above.

### 13.4 Multicollinearity*

## Reference: Greene (2000) 6.7

When the variables in the $x_{t}$ vector are very highly correlated (they are "multicollinear") then data cannot tell, with the desired precision, if the movements in $y_{t}$ was due to movements in $x_{i t}$ or $x_{j t}$. This means that the point estimates might fluctuate wildly over subsamples and it is often the case that individual coefficients are insignificant even though the $R^{2}$ is high and the joint significance of the coefficients is also high. The estimators are still consistent and asymptotically normally distributed, just very imprecise.

A common indicator for multicollinearity is to standardize each element in $x_{t}$ by subtracting the sample mean and then dividing by its standard deviation

$$
\begin{equation*}
\tilde{x}_{i t}=\left(x_{i t}-\bar{x}_{i t}\right) / \operatorname{std}\left(x_{i t}\right) . \tag{13.13}
\end{equation*}
$$

(Another common procedure is to use $\tilde{x}_{i t}=x_{i t} /\left(\Sigma_{t=1}^{T} x_{i t}^{2} / T\right)^{1 / 2}$.)
Then calculate the eigenvalues, $\lambda_{j}$, of the second moment matrix of $\tilde{x}_{t}$

$$
\begin{equation*}
A=\frac{1}{T} \sum_{t=1}^{T} \tilde{x}_{t} \tilde{x}_{t}^{\prime} \tag{13.14}
\end{equation*}
$$

The condition number of a matrix is the ratio of the largest (in magnitude) of the eigenvalues to the smallest

$$
\begin{equation*}
c=|\lambda|_{\max } /|\lambda|_{\min } \tag{13.15}
\end{equation*}
$$

(Some authors take $c^{1 / 2}$ to be the condition number; others still define it in terms of the "singular values" of a matrix.) If the regressors are uncorrelated, then the condition value of $A$ is one. This follows from the fact that $A$ is a (sample) covariance matrix. If it is diagonal, then the eigenvalues are equal to diagonal elements, which are all unity since the standardization in (13.13) makes all variables have unit variances. Values of $c$ above several hundreds typically indicate serious problems.

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## 14 Generalized Least Squares

Reference: Greene (2000) 11.3-4
Additional references: Hayashi (2000) 1.6; Johnston and DiNardo (1997) 5.4; Verbeek (2000) 6

### 14.1 Introduction

Instead of using LS in the presence of autocorrelation/heteroskedasticity (and, of course, adjusting the variance-covariance matrix), we may apply the generalized least squares method. It can often improve efficiency.

The linear model $y_{t}=x_{t}^{\prime} \beta_{0}+u_{t}$ written on matrix form (GLS is one of the cases in econometrics where matrix notation really pays off) is

$$
\begin{aligned}
& y=X \beta_{0}+u, \text { where } \\
& y=\left[\begin{array}{l}
y_{1} \\
y_{2} \\
\vdots \\
y_{T}
\end{array}\right], X=\left[\begin{array}{l}
x_{1}^{\prime} \\
x_{2}^{\prime} \\
\vdots \\
x_{T}^{\prime}
\end{array}\right], \text { and } u=\left[\begin{array}{l}
u_{1} \\
u_{2} \\
\vdots \\
u_{T}
\end{array}\right]
\end{aligned}
$$

Suppose that the covariance matrix of the residuals (across time) is

$$
\begin{align*}
\mathrm{E} u u^{\prime} & =\left[\begin{array}{cccc}
\mathrm{E} u_{1} u_{1} & \mathrm{E} u_{1} u_{2} & \cdots & \mathrm{E} u_{1} u_{T} \\
\mathrm{E} u_{2} u_{1} & \mathrm{E} u_{2} u_{2} & & \mathrm{E} u_{2} u_{T} \\
\vdots & & \ddots & \vdots \\
\mathrm{E} u_{T} u_{1} & \mathrm{E} u_{T} u_{2} & & \mathrm{E} u_{T} u_{T}
\end{array}\right] \\
& =\Omega_{T \times T .} . \tag{14.2}
\end{align*}
$$

This allows for both heteroskedasticity (different elements along the main diagonal) and autocorrelation (non-zero off-diagonal elements). LS is still consistent even if $\Omega$ is not proportional to an identity matrix, but it is not efficient. Generalized least squares (GLS)
is. The trick of GLS is to transform the variables and the do LS.

### 14.2 GLS as Maximum Likelihood

Remark 1 If the $n \times 1$ vector $x$ has a multivariate normal distribution with mean vector $\mu$ and covariance matrix $\Omega$, then the joint probability density function is $(2 \pi)^{-n / 2}|\Omega|^{-1 / 2} \exp [-(x-$ $\left.\mu)^{\prime} \Omega^{-1}(x-\mu) / 2\right]$.

If the $T \times 1$ vector $u$ is $N(\mathbf{0}, \Omega)$, then the joint pdf of $u$ is $(2 \pi)^{-n / 2}|\Omega|^{-1 / 2} \exp \left[-u^{\prime} \Omega^{-1} u / 2\right]$.
Change variable from $u$ to $y-X \beta$ (the Jacobian of this transformation equals one), and take logs to get the (scalar) log likelihood function

$$
\begin{equation*}
\ln L=-\frac{n}{2} \ln (2 \pi)-\frac{1}{2} \ln |\Omega|-\frac{1}{2}(y-X \beta)^{\prime} \Omega^{-1}(y-X \beta) . \tag{14.3}
\end{equation*}
$$

To simplify things, suppose we know $\Omega$. It is then clear that we maximize the likelihood function by minimizing the last term, which is a weighted sum of squared errors.

In the classical $L S$ case, $\Omega=\sigma^{2} I$, so the last term in (14.3) is proportional to the unweighted sum of squared errors. The LS is therefore the MLE when the errors are iid normally distributed.

When errors are heteroskedastic, but not autocorrelated, then $\Omega$ has the form

$$
\Omega=\left[\begin{array}{cccc}
\sigma_{1}^{2} & 0 & \cdots & 0  \tag{14.4}\\
0 & \sigma_{2}^{2} & & \vdots \\
\vdots & & \ddots & 0 \\
0 & \cdots & 0 & \sigma_{T}^{2}
\end{array}\right]
$$

In this case, we can decompose $\Omega^{-1}$ as

$$
\Omega^{-1}=P^{\prime} P \text {, where } P=\left[\begin{array}{cccc}
1 / \sigma_{1} & 0 & \cdots & 0 \\
0 & 1 / \sigma_{2} & & \vdots \\
\vdots & & \ddots & 0 \\
0 & \cdots & 0 & 1 / \sigma_{T}
\end{array}\right]
$$

The last term in (14.3) can then be written

$$
\begin{align*}
-\frac{1}{2}(y-X \beta)^{\prime} \Omega^{-1}(y-X \beta) & =-\frac{1}{2}(y-X \beta)^{\prime} P^{\prime} P(y-X \beta) \\
& =-\frac{1}{2}(P y-P X \beta)^{\prime}(P y-P X \beta) \tag{14.6}
\end{align*}
$$

This very practical result says that if we define $y_{t}^{*}=y_{t} / \sigma_{t}$ and $x_{t}^{*}=x_{t} / \sigma_{t}$, then we get ML estimates of $\beta$ running an LS regression of $y_{t}^{*}$ on $x_{t}^{*}$. (One of the elements in $x_{t}$ could be a constant-also this one should be transformed). This is the generalized least squares (GLS).

Remark 2 Let A be an $n \times n$ symmetric positive definite matrix. It can be decomposed as $A=P P^{\prime}$. There are many such $P$ matrices, but only one which is lower triangular $P$ (see next remark)

Remark 3 Let A be an $n \times n$ symmetric positive definite matrix. The Cholesky decomposition gives the unique lower triangular $P_{1}$ such that $A=P_{1} P_{1}^{\prime}$ or an upper triangular matrix $P_{2}$ such that $A=P_{2}^{\prime} P_{2}$ (clearly $P_{2}=P_{1}^{\prime}$ ). Note that $P_{1}$ and $P_{2}$ must be invertible (since $A$ is).

When errors are autocorrelated (with or without heteroskedasticity), then it is typically harder to find a straightforward analytical decomposition of $\Omega^{-1}$. We therefore move directly to the general case. Since the covariance matrix is symmetric and positive definite, $\Omega^{-1}$ is too. We therefore decompose it as

$$
\begin{equation*}
\Omega^{-1}=P^{\prime} P . \tag{14.7}
\end{equation*}
$$

The Cholesky decomposition is often a convenient tool, but other decompositions can also be used. We can then apply (14.6) also in this case-the only difference is that $P$ is typically more complicated than in the case without autocorrelation. In particular, the transformed variables $P y$ and $P X$ cannot be done line by line ( $y_{t}^{*}$ is a function of $y_{t}, y_{t-1}$, and perhaps more).

Example 4 (AR(1) errors, see Davidson and MacKinnon (1993) 10.6.) Let $u_{t}=a u_{t-1}+$ $\varepsilon_{t}$ where $\varepsilon_{t}$ is iid. We have $\operatorname{Var}\left(u_{t}\right)=\sigma^{2} /\left(1-a^{2}\right)$, and $\operatorname{Corr}\left(u_{t}, u_{t-s}\right)=a^{s}$. For $T=4$,
the covariance matrix of the errors is

$$
\begin{aligned}
\Omega & =\operatorname{Cov}\left(\left[\begin{array}{llll}
u_{1} & u_{2} & u_{3} & u_{4}
\end{array}\right]^{\prime}\right) \\
& =\frac{\sigma^{2}}{1-a^{2}}\left[\begin{array}{cccc}
1 & a & a^{2} & a^{3} \\
a & 1 & a & a^{2} \\
a^{2} & a & 1 & a \\
a^{3} & a^{2} & a & 1
\end{array}\right]
\end{aligned}
$$

The inverse is

$$
\Omega^{-1}=\frac{1}{\sigma^{2}}\left[\begin{array}{cccc}
1 & -a & 0 & 0 \\
-a & 1+a^{2} & -a & 0 \\
0 & -a & 1+a^{2} & -a \\
0 & 0 & -a & 1
\end{array}\right]
$$

and note that we can decompose it as

This is not a Cholesky decomposition, but certainly a valid decomposition (in case of doubt, do the multiplication). Premultiply the system

$$
\left[\begin{array}{l}
y_{1} \\
y_{2} \\
y_{3} \\
y_{4}
\end{array}\right]=\left[\begin{array}{l}
x_{1}^{\prime} \\
x_{2}^{\prime} \\
x_{3}^{\prime} \\
x_{4}^{\prime}
\end{array}\right] \beta_{0}+\left[\begin{array}{l}
u_{1} \\
u_{2} \\
u_{3} \\
u_{4}
\end{array}\right]
$$

by P to get

$$
\frac{1}{\sigma}\left[\begin{array}{c}
\sqrt{\left(1-a^{2}\right)} y_{1} \\
y_{2}-a y_{1} \\
y_{3}-a y_{2} \\
y_{4}-a y_{3}
\end{array}\right]=\frac{1}{\sigma}\left[\begin{array}{c}
\sqrt{\left(1-a^{2}\right)} x_{1}^{\prime} \\
x_{2}^{\prime}-a x_{1}^{\prime} \\
x_{3}^{\prime}-a x_{2}^{\prime} \\
x_{4}^{\prime}-a x_{3}^{\prime}
\end{array}\right] \beta_{0}+\frac{1}{\sigma}\left[\begin{array}{c}
\sqrt{\left(1-a^{2}\right)} u_{1} \\
\varepsilon_{2} \\
\varepsilon_{3} \\
\varepsilon_{4}
\end{array}\right]
$$

Note that all the residuals are uncorrelated in this formulation. Apart from the first observation, they are also identically distributed. The importance of the first observation becomes smaller as the sample size increases-in the limit, GLS is efficient.

### 14.3 GLS as a Transformed LS

When the errors are not normally distributed, then the MLE approach in the previous section is not valid. But we can still note that GLS has the same properties as LS has with iid non-normally distributed errors. In particular, the Gauss-Markov theorem applies, so the GLS is most efficient within the class of linear (in $y_{t}$ ) and unbiased estimator (assuming, of course, that GLS and LS really are unbiased, which typically requires that $u_{t}$ is uncorrelated with $x_{t-s}$ for all $s$ ). This follows from that the transformed system

$$
\begin{align*}
P y & =P X \beta_{0}+P u \\
y^{*} & =X^{*} \beta_{0}+u^{*}, \tag{14.8}
\end{align*}
$$

have iid errors, $u^{*}$. So see this, note that

$$
\begin{align*}
\mathrm{E} u^{*} u^{* \prime} & =\mathrm{E} P u u^{\prime} P^{\prime} \\
& =P \mathrm{E} u u^{\prime} P^{\prime} \tag{14.9}
\end{align*}
$$

Recall that $\mathrm{E} u u^{\prime}=\Omega, P^{\prime} P=\Omega^{-1}$ and that $P^{\prime}$ is invertible. Multiply both sides by $P^{\prime}$

$$
\begin{align*}
P^{\prime} \mathrm{E} u^{*} u^{* \prime} & =P^{\prime} P \mathrm{E} u u^{\prime} P^{\prime} \\
& =\Omega^{-1} \Omega P^{\prime} \\
& =P^{\prime}, \text { so } \mathrm{E} u^{*} u^{* \prime}=I . \tag{14.10}
\end{align*}
$$

### 14.4 Feasible GLS

In practice, we usually do not know $\Omega$. Feasible $G L S$ (FGSL) is typically implemented by first estimating the model (14.1) with LS, then calculating a consistent estimate of $\Omega$, and finally using GLS as if $\Omega$ was known with certainty. Very little is known about the finite sample properties of FGLS, but (the large sample properties) consistency, asymptotic normality, and asymptotic efficiency (assuming normally distributed errors) can often be
established. Evidence from simulations suggests that the FGLS estimator can be a lot worse than LS if the estimate of $\Omega$ is bad.

To use maximum likelihood when $\Omega$ is unknown requires that we make assumptions about the structure of $\Omega$ (in terms of a small number of parameters), and more generally about the distribution of the residuals. We must typically use numerical methods to maximize the likelihood function

Example 5 (MLE and AR(1) errors.) If $u_{t}$ in Example 4 are normally distributed, then we can use the $\Omega^{-1}$ in (14.3) to express the likelihood function in terms of the unknown parameters: $\beta, \sigma$, and $a$. Maximizing this likelihood function requires a numerical optimization routine.

## Bibliography

Davidson, R., and J. G. MacKinnon, 1993, Estimation and Inference in Econometrics, Oxford University Press, Oxford

Greene, W. H., 2000, Econometric Analysis, Prentice-Hall, Upper Saddle River, New Jersey, 4th edn.

## Hayashi, F., 2000, Econometrics, Princeton University Press.

Johnston, J., and J. DiNardo, 1997, Econometric Methods, McGraw-Hill, New York, 4th edn.

Verbeek, M., 2000, A Guide to Modern Econometrics, Wiley, Chichester.

## 0 Reading List

Main reference: Greene (2000) (GR)
$\left(^{*}\right)$ denotes required reading.

### 0.1 Introduction

1. *Lecture notes

### 0.2 Time Series Analysis

1. *Lecture notes
2. *GR 13.1-13.3, 18.1-18.2, 17.5
3. Obstfeldt and Rogoff (1996) 2.3.5
4. $\operatorname{Sims}(1980)$

Keywords: moments of a time series process, covariance stationarity, ergodicity, conditional and unconditional distributions, white noise, MA, AR, MLE of AR process, VAR (Advanced: unit roots, cointegration)

### 0.3 Distribution of Sample Averages

1. ${ }^{*}$ Lecture notes
2. GR 11.2

Keywords: Newey-West

### 0.4 Asymptotic Properties of LS

1. *Lecture notes
2. *GR 9.1-9.4, 11.2

Keywords: consistency of LS, asymptotic normality of LS, influential observations, robust estimators, LAD

### 0.5 Instrumental Variable Method

1. *Lecture notes
2. *GR 9.5 and 16.1-2

Keywords: measurement errors, simultaneous equations bias, instrumental variables, 2SLS
0.6 Simulating the Finite Sample Properties

1. *Lecture notes
2. ${ }^{*}$ GR 5.3

Keywords: Monte Carlo simulations, Bootstrap simulations
0.7 GMM

1. *Lecture notes
2. *GR 4.7 and 11.5-6
3. Christiano and Eichenbaum (1992)

Keywords: method of moments, unconditional/conditional moment conditions, loss function, asymptotic distribution of GMM estimator, efficient GMM, GMM and inference
0.7.1 Application of GMM: LS/IV with Autocorrelation and Heteroskedasticity

1. *Lecture notes
2. ${ }^{*}$ GR 12.2 and 13.4
3. Lafontaine and White (1986)
4. Mankiw, Romer, and Weil (1992)

Keywords: finite sample properties of LS and IV, consistency of LS and IV, asymptotic distribution of LS and IV

### 0.7.2 Application of GMM: Systems of Simultaneous Equations

1. *Lecture notes
2. *GR 16.1-2, 16.3 (introduction only)
3. Obstfeldt and Rogoff (1996) 2.1
4. Deaton (1992) 3

Keywords: structural and reduced forms, identification, 2SLS

## Bibliography

Christiano, L. J., and M. Eichenbaum, 1992, "Current Real-Business-Cycle Theories and Aggregate Labor-Market Fluctuations," American Economic Review, 82, 430-450.

Deaton, A., 1992, Understanding Consumption, Oxford University Press.
Greene, W. H., 2000, Econometric Analysis, Prentice-Hall, Upper Saddle River, New Jersey, 4th edn.

Lafontaine, F., and K. J. White, 1986, "Obtaining Any Wald Statistic You Want," Economics Letters, 21, 35-40.

Mankiw, N. G., D. Romer, and D. N. Weil, 1992, "A Contribution to the Empirics of Economic Growth," Quarterly Journal of Economics, 107, 407-437.

Obstfeldt, M., and K. Rogoff, 1996, Foundations of International Macroeconomics, MIT Press.

Sims, C. A., 1980, "Macroeconomics and Reality," Econometrica, 48, 1-48.


[^0]:    ${ }^{1}$ The relation between $f_{t}\left(y_{t}\right)$ and $g_{t-s, t}\left(y_{t-s}, y_{t}\right)$ is, as usual, $f_{t}\left(y_{t}\right)=\int_{-\infty}^{\infty} g_{t-s, t}\left(y_{t-s}, y_{t}\right) d y_{t-s}$

[^1]:    ${ }^{1}$ This puts non-trivial restrictions on the data generating processes. For instance, if $x_{t}$ include lagged values of $y_{t}$, then we typically require $y_{t}$ to be stationary and ergodic, and that $u_{t}$ is independent of $x_{t-s}$ for $s \geq 0$.

[^2]:    ${ }^{2}$ However, some simulation evidence suggests that F tests may have better small sample properties than
    chi-square test. chi-square test.

[^3]:    ${ }^{1}$ In matrix notation where $z_{t}^{\prime}$ is the $t^{\text {th }}$ row of $Z$ we have $\hat{\beta}_{I V}=\left(Z^{\prime} X / T\right)^{-1}\left(Z^{\prime} Y / T\right)$.

[^4]:    ${ }^{1}$ By premultiplying with $F^{-1}$ and rearranging we get the reduced form $y_{t}=\Pi z_{t}+\varepsilon_{t}$, with $\Pi=-F^{-1}$

[^5]:    ${ }^{1}$ This is a "structural model" in a traditional, Cowles commission, sense. This might be different from what modern macroeconomists would call structural.

