

25 Spring 439/639 TSA: Lecture 11

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More transformations of time series

The goal of transformation is to make the time series “more” stationary, and/or make the time series a normal process (if possible). Last time we already introduced **Variance Stabilizing Transformation** (under certain settings).

Another option: **Box–Cox transformations**, transforming y to $g(y)$ as follows

$$g(y) = \begin{cases} \frac{y^\lambda - 1}{\lambda}, & \lambda \neq 0 \\ \log y, & \lambda = 0 \end{cases}$$

Exercise: Show $\lim_{\lambda \rightarrow 0} \frac{y^\lambda - 1}{\lambda} = \log y$.

The λ in Box–Cox transformation above is chosen via an MLE approach. See R notebook later.

Another common way is to **take the difference of logarithm**, i.e., log-differences. This can be practically useful in specific applications.

In finance, suppose the time series (Y_t) can be written as follows

$$Y_t = Y_{t-1} + X_t \cdot Y_{t-1} = Y_{t-1} (1 + X_t).$$

So (X_t) is the percentage change of (Y_t) , and we have

$$\log Y_t = \log Y_{t-1} + \log(1 + X_t).$$

Then the log-difference, (or the log-returns, the returns) of (Y_t) is

$$\nabla \log Y_t = \log Y_t - \log Y_{t-1} = \log(1 + X_t) \approx X_t$$

where the last step is because the percentage change X_t is usually small in finance. (*Note:* X_t can be positive or negative, but close to 0.) And in practice, the time series (X_t) is usually stationary. So taking the transformation $\nabla \log Y_t$ gives a more stationary time series.

Summary: to make a time series more stationary, we can consider differencing, variance stabilizing transformation, taking logarithm, Box-Cox transformations, and combination of these.

Roadmap ahead

So far in this course, we already studied models for time series (we mainly studied $\text{ARIMA}(p, d, q)$).

In the next module of this course, we will look at the following topics: Suppose we observe (y_1, y_2, \dots, y_n) ,

1. **Model Specification:** determine the best/ the most appropriate (p, d, q) in the model, based on sample ACF or sample PACF.
2. **Fit the selected model or Parameter Estimation:** using MLE, LS, MoM, etc.
3. **Diagnostic:** analyze the residuals.
4. **Forecasting**

Bartlett's Theorem

First, recall the sample ACF for observed (Y_1, \dots, Y_n) : (see lecture 9)

$$\hat{\rho}_k = r_k = \frac{\sum_{t=k+1}^n (Y_t - \bar{Y})(Y_{t-k} - \bar{Y})}{\sum_{t=1}^n (Y_t - \bar{Y})^2}.$$

Remark: The theoretical ACF ρ_k of a given time series model is a non-random number (but it can be unknown). The sample ACF r_k , defined above, is a random variable if we think (Y_1, \dots, Y_n) as a random realization of a time series model. So the sample ACF r_k follows a certain sampling distribution (which depends on the underlying time series model).

The **Bartlett's Theorem** (see equation (6.1.2) in Cryan and Chan) says, for a fixed m , the sampling (joint) distribution of (r_1, \dots, r_m) approaches a multivariate normal distribution as $n \rightarrow \infty$, i.e.,

$$\vec{r} \sim \text{MVN}(\vec{\rho}, \frac{1}{n}C), \quad \text{as } n \rightarrow \infty,$$

$$\text{where } \vec{r} = \begin{bmatrix} r_1 \\ r_2 \\ \vdots \\ r_m \end{bmatrix}, \quad \vec{\rho} = \begin{bmatrix} \rho_1 \\ \rho_2 \\ \vdots \\ \rho_m \end{bmatrix},$$

n is the sample size of observations (of the time series), and the matrix C is an $m \times m$ matrix. The full detailed formula for C can be found in equation (6.1.2) in Cryan and Chan. (*Note:* In general, C is not diagonal.)

In particular, the diagonal entries of C are

$$c_{ii} = \sum_{k=-\infty}^{+\infty} (\rho_{k+i}^2 + \rho_{k-i}\rho_{k+i} - 4\rho_i\rho_k\rho_{k+i} + 2\rho_i^2\rho_k^2).$$

Then the Bartlett's Theorem implies that, as $n \rightarrow \infty$,

$$r_i \sim \mathcal{N}\left(\rho_i, \frac{1}{n}c_{ii}\right).$$

Example: white noise

Suppose $Y_t \sim \text{WN}(0, \sigma^2)$, then

$$\rho_0 = 1, \text{ and } \rho_i = 0 \text{ for } i \geq 1.$$

Let's look at

$$c_{ii} = \sum_{k=-\infty}^{+\infty} (\rho_{k+i}^2 + \rho_{k-i}\rho_{k+i} - 4\rho_i\rho_k\rho_{k+i} + 2\rho_i^2\rho_k^2).$$

If $i \geq 1$, then $\rho_i = 0$, $c_{ii} = \sum_{k=-\infty}^{+\infty} (\rho_{k+i}^2 + \rho_{k-i}\rho_{k+i})$. Also note that $\rho_{k-i}\rho_{k+i} = 0$ for any k , and $\sum_{k=-\infty}^{+\infty} \rho_{k+i}^2 = \rho_0^2 = 1$, so $c_{ii} = 1$.

By Bartlett's Theorem, for any fixed $i \geq 1$, $r_i \sim N(0, \frac{1}{n})$ (when sample size n is large). Using this result, we can construct 95% CI for ρ_i (for this example):

$$\left[r_i - \frac{2}{\sqrt{n}}, r_i + \frac{2}{\sqrt{n}} \right].$$

We can also compute c_{00} . If $i = 0$, then $c_{ii} = \sum_{k=-\infty}^{+\infty} (\rho_{k+i}^2 + \rho_k^2 - 4\rho_k^2 + 2\rho_k^2) = 0$. By a similar statement from Bartlett's Theorem, $\text{Var}(r_0) = 0$. This is not a surprise since the sample ACF at lag 0, i.e. r_0 , is always 1.

Example: AR(1)

Suppose (Y_t) follows AR(1), then

$$\rho_k = \phi^k \text{ for } k \geq 0, \text{ and } \rho_k = \phi^{|k|} \text{ for } k < 0.$$

Using the c_{ii} formula from Bartlett's Theorem, we can derive that: (for large n)

$$\text{Var}(r_i) = \frac{c_{ii}}{n} = \frac{1}{n} \left[\frac{(1 + \phi^2)(1 - \phi^{2i})}{1 - \phi^2} - 2i\phi^{2i} \right].$$

In particular, when $i = 1$, we have $\text{Var}(r_1) = \frac{1 - \phi^2}{n}$. So if ϕ is close to ± 1 , then r_1 is a very precise estimate of ρ_1 .

If the lag i is very large, then $\text{Var}(r_i) \approx \frac{1}{n} \cdot \frac{1 + \phi^2}{1 - \phi^2}$. So if ϕ is close to ± 1 , then for large i , r_i is not a precise estimate of ρ_i (in the sense that the variance is very large).

Example: MA(1)

Suppose (Y_t) follows MA(1), then

$$\rho_0 = 1, \quad \rho_1 = \rho_{-1} \neq 0, \quad \text{and } \rho_k = 0 \text{ for } |k| \geq 2.$$

For c_{11} , we have

$$\begin{aligned} c_{11} &= \sum_{k=-\infty}^{+\infty} (\rho_{k+1}^2 + \rho_{k-1}\rho_{k+1} - 4\rho_1\rho_k\rho_{k+1} + 2\rho_1^2\rho_k^2) \\ &= (\rho_0^2 + 2\rho_1^2) + (\rho_{-1}\rho_1) - 4\rho_1(\rho_{-1}\rho_0 + \rho_0\rho_1) + 2\rho_1^2(\rho_0^2 + 2\rho_1^2) \\ &= 1 + 2\rho_1^2 + \rho_1^2 - 4\rho_1^2 + 2\rho_1^2 + 4\rho_1^4 \\ &= 1 - 3\rho_1^2 + 4\rho_1^4. \end{aligned}$$

By Bartlett's Theorem, we have the following (for large n)

$$r_1 \sim \mathcal{N}\left(\rho_1, \frac{1 - 3\rho_1^2 + 4\rho_1^4}{n}\right).$$

For $i \geq 2$, we can also derive that

$$c_{ii} = 1 + 2\rho_1^2, \quad r_i \sim \mathcal{N}\left(0, \frac{1 + 2\rho_1^2}{n}\right).$$

Exercise: verify that $c_{ii} = 1 + 2\rho_1^2$ for any $i \geq 2$ (under the MA(1) setting).

Example: MA(q)

For MA(q), we can show that: (for large n)

$$r_i \sim \mathcal{N}\left(0, \frac{1 + 2\sum_{j=1}^q \rho_j^2}{n}\right), \text{ for any lag } i \geq q + 1,$$

which is similar to the $i \geq 2$ case in MA(1).

Hypothesis Testing for MA(q)

We want to do hypothesis testing in the following form

$$H_0 : \text{series is MA}(q) \quad \text{vs.} \quad H_a : \text{series is not MA}(q).$$

For example, let's first look at the hypothesis testing for MA(1):

$$H_0 : \text{series is MA}(1) \quad \text{vs.} \quad H_a : \text{series is not MA}(1).$$

By the earlier results from this lecture, under H_0 , i.e. MA(1),

$$r_i \sim \mathcal{N}\left(0, \frac{1 + 2\rho_1^2}{n}\right), \quad \text{for any } i \geq 2.$$

So under H_0 , with 95% probability,

$$r_i \in \left[-\frac{2}{\sqrt{n}}\sqrt{1 + 2\rho_1^2}, \frac{2}{\sqrt{n}}\sqrt{1 + 2\rho_1^2}\right],$$

and approximately, with 95% probability,

$$r_i \in \left[-\frac{2}{\sqrt{n}}\sqrt{1 + 2r_1^2}, \frac{2}{\sqrt{n}}\sqrt{1 + 2r_1^2}\right].$$

So we construct the rejection region

$$|r_i| > \frac{2}{\sqrt{n}}\sqrt{1 + 2r_1^2}, \text{ for } i \geq 2.$$

For example if $|r_2| > \frac{2}{\sqrt{n}}\sqrt{1 + 2r_1^2}$, then reject H_0 . (*Note:* we may have access to multiple sample ACFs (r_2, r_3, \dots) , rather than just one single r_2 . We can use them collectively to reject H_0 .)

If the H_0 of the previous hypothesis testing, i.e. MA(1) is rejected, then we turn to the next hypothesis testing:

$$H_0 : \text{series is MA}(2) \quad \text{vs.} \quad H_a : \text{series is not MA}(2).$$

We can use a similar analysis to derive the rejection region (see the earlier example of MA(q) in this lecture) for this hypothesis testing. The rejection region for H_0 (series is MA(2)) is

$$|r_i| > \frac{2}{\sqrt{n}}\sqrt{1 + 2r_1^2 + 2r_2^2}, \text{ for } i \geq 3.$$

We can repeat this process until we fail to reject a hypothesis that the series is an $MA(q)$ (for some q).

Example: Assume we observe a sample with $n = 100$, and the sample ACFs are $r_1 = 0.5$, $r_2 = 0.4$, $r_3 = 0.4$, $r_4 = 0.3$. Suppose we want to hypothesis testing to determine an $MA(q)$ model for this sample.

We start with testing whether it is white noise (i.e. $MA(0)$). For white noise, $r_i \in [-\frac{2}{\sqrt{n}}, \frac{2}{\sqrt{n}}]$ for all $i \geq 1$ with 95% probability. (The rejection region for white noise is $|r_i| > \frac{2}{\sqrt{n}} = 0.2$ for $i \geq 1$.) We reject the white noise using the observed sample ACFs r_1, r_2, r_3, r_4 .

For testing $MA(1)$: the rejection region is

$$|r_i| > \frac{2}{\sqrt{n}} \sqrt{1 + 2r_1^2} = \frac{2}{10} \sqrt{1 + 2 \cdot 0.5^2} = 0.2\sqrt{1.5} \approx 0.245, \quad \text{for } i \geq 2.$$

So we reject the $MA(1)$ using the observed sample ACFs r_2, r_3, r_4 . Next, test $MA(2)$.

Exercise: finish this testing process for this example.

Brief introduction to Partial Autocorrelation Function (PACF)

Motivation: We just saw that one can test if a process is an $MA(q)$ for a specific q using the sample ACFs. A natural question is, can we determine the order p of $AR(p)$ from the sample ACFs? The answer is no. And we need to look at Partial ACF (PACF) to determine the order p of $AR(p)$. The sample PACF can help us determine the order of $AR(p)$.

The notation of **Partial Autocorrelation Function (PACF)** is ϕ_{kk} , which denotes the partial autocorrelation at lag k . We will see several different definitions in the next lecture. Here is one way to define it.

$$\phi_{kk} \stackrel{\text{def}}{=} \text{corr}(Y_t, Y_{t-k} \mid Y_{t-1}, Y_{t-2}, \dots, Y_{t-k+1}).$$

This means, ϕ_{kk} is the **conditional correlation** between Y_t and Y_{t-k} **conditional on** all intermediate values $Y_{t-1}, Y_{t-2}, \dots, Y_{t-k+1}$.

Example: consider the $AR(1)$ process, $Y_t = \phi Y_{t-1} + e_t$ (with $\phi \neq 0$). We can show that $\phi_{11} = \rho_1 \neq 0$, and $\phi_{kk} = 0$ for any lag $k \geq 2$.

In general, for $AR(p)$, the PACF at lags 1 through p can be nonzero, and PACF at lags $k \geq p + 1$ are all zero.

Remark: for any time series (not necessarily special models like ARIMA), the definition itself contains two special cases. For $k = 0$, ϕ_{00} is always 1 by definition. For $k = 1$, the conditional correlation reduces to an unconditional correlation, so $\phi_{11} = \text{corr}(Y_t, Y_{t-1})$, and $\phi_{11} = \rho_1$ assuming stationarity.