# CHAPTER 4. LIMIT THEOREMS IN STATISTICS 

### 4.1. SEQUENCES OF RANDOM VARIABLES

4.1.1. A great deal of econometrics uses relatively large data sets and methods of statistical inference that are justified by their desirable properties in large samples. The probabilistic foundations for these arguments are "laws of large numbers", sometimes called the "law of averages", and "central limit theorems". This chapter presents these foundations. It concentrates on the simplest versions of these results, but goes some way in covering more complicated versions that are needed for some econometric applications. For basic econometrics, the most critical materials are the limit concepts and their relationship covered in this section, and for independent and identically distributed (i.i.d.) random variables the first Weak Law of Large Numbers in Section 4.3 and the first Central Limit Theorem in Section 4.4. The reader may want to postpone other topics, and return to them as they are needed in later chapters.
4.1.2. Consider a sequence of random variables $\mathrm{Y}_{1}, \mathrm{Y}_{2}, \mathrm{Y}_{3}, \ldots$. These random variables are all functions $\mathrm{Y}_{\mathrm{k}}(\mathrm{s})$ of the same state of Nature s , but may depend on different parts of s . There are several possible concepts for the limit $Y_{o}$ of a sequence of random variables $Y_{n}$. Since the $Y_{n}$ are functions of states of nature, these limit concepts will correspond to different ways of defining limits of functions. Figure 4.1 will be used to discuss limit concepts. Panel (a) graphs $Y_{n}$ and $Y_{o}$ as functions of the state of Nature. Also graphed are curves denoted $Y_{o \pm}$ and defined by $Y_{o} \pm \varepsilon$ which for each state of Nature s delineate an $\varepsilon$-neighborhood of $Y_{o}(s)$. The set of states of Nature for which $\left|Y_{o}(s)-Y_{n}(s)\right|>\varepsilon$ is denoted $\mathbf{W}_{n}$. Panel (b) graphs the CDF's of $Y_{o}$ and $Y_{n}$. For technical completeness, note that a random variable Y is a measurable real-valued function on a probability space (S,F,P), where $\boldsymbol{F}$ is a $\boldsymbol{\sigma}$-field of subsets of $\mathbf{S}, \mathrm{P}$ is a probability on $\boldsymbol{F}$, and "measurable" means that F contains the inverse image of every set in the Borel $\sigma$-field of subsets of the real line. The CDF of a vector of random variables is then a measurable function with the properties given in 3.5.3.
4.1.3. $\mathrm{Y}_{\mathrm{n}}$ converges in probability to $\mathrm{Y}_{\mathrm{o}}$, if for each $\varepsilon>0, \lim _{\mathrm{n} \rightarrow \infty} \operatorname{Prob}\left(\left|\mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}}\right|>\varepsilon\right)=0$. Convergence in probability is denoted $Y_{n} \rightarrow_{p} Y_{o}$, or $\operatorname{plim}_{n \rightarrow \infty} Y_{n}=Y_{o}$. With $W_{n}$ defined as in Figure 4.1, $\mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{p}} \mathrm{Y}_{\mathrm{o}}$ iff $\lim _{\mathrm{n} \rightarrow \infty} \operatorname{Prob}\left(\mathbf{W}_{\mathrm{n}}\right)=0$ for each $\varepsilon>0$.
4.1.4. $\mathrm{Y}_{\mathrm{n}}$ converges almost surely to $\mathrm{Y}_{\mathrm{o}}$, denoted $\mathrm{Y}_{\mathrm{n}} \rightarrow_{\text {as }} \mathrm{Y}_{0}$, if for each $\varepsilon>0$, $\lim _{\mathrm{n} \rightarrow \infty} \operatorname{Prob}\left(\sup _{\mathrm{m}>\mathrm{n}}\left|\mathrm{Y}_{\mathrm{m}}-\mathrm{Y}_{\mathrm{o}}\right|>\varepsilon\right)=0$. For $\mathbf{W}_{\mathrm{n}}$ defined in Figure 4.1, the set of states of nature for which $\left|\mathrm{Y}_{\mathrm{m}}(\mathrm{w})-\mathrm{Y}_{\mathrm{o}}(\mathrm{w})\right|>\varepsilon$ for some $\mathrm{m} \geq \mathrm{n}$ is $\mathrm{U}_{m \geq n} \mathbf{W}_{\mathrm{m}}$, and $\mathrm{Y}_{\mathrm{n}} \rightarrow_{\text {as }} \mathrm{Y}_{\mathrm{o}}$ iff $\operatorname{Prob}\left(U_{m \geq n} \mathbf{W}_{\mathrm{n}^{\prime}}\right) \rightarrow 0$.

An implication of almost sure convergence is $\lim _{n \rightarrow \infty} Y_{n}(s)=Y_{o}(s)$ a.s. (i.e., except for a set of states of Nature of probability zero); this is not an implication of $Y_{n}{ }_{p} Y_{o}$.

## FIGURE 4.1. CONVERGENCE CONCEPTS FOR RANDOM VARIABLES



Panel (a)

Panel (b)

4.1.5. $\mathrm{Y}_{\mathrm{n}}$ converges in $\rho$-mean (also called convergence in $\|\cdot\|_{\rho}$ norm, or convergence in $\mathbf{L}_{\mathrm{p}}$ space) to $\mathrm{Y}_{\mathrm{o}}$ if $\lim _{\mathrm{n} \rightarrow \infty} \mathbf{E}\left|\mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}}\right|^{\rho}=0$. For $\rho=2$, this is called convergence in quadratic mean. The norm is defined as $\|\mathrm{Y}\|_{\rho}=\left[\int_{S}|\mathrm{Y}(\mathrm{s})|^{\rho \cdot} \cdot \mathrm{P}(\mathrm{ds})\right]^{1 / \rho}=\left[\mathbf{E}|\mathrm{Y}|^{\rho}\right]^{1 / \mathrm{\rho}}$, and can be interpreted as a probabilityweighted measure of the distance of Y from zero. The norm of a random variable is a moment. There are random variables for which the $\rho$-mean will not exist for any $\rho>0$; for example, Y with $\operatorname{CDF} F(\mathrm{y})=1-1 /(\log \mathrm{y})$ for $\mathrm{y} \geq e$ has this property. However, in many applications moments such as variances exist, and the quadratic mean is a useful measure of distance.
4.1.6. $Y_{n}$ converges in distribution to $Y_{o}$, denoted $Y_{n}{ }_{d} Y_{o}$, if the CDF of $Y_{n}$ converges to the CDF of $Y_{o}$ at each continuity point of $Y_{o}$. In Figure $4.1(b)$, this means that $F_{n}$ converges to the function $\mathrm{F}_{\mathrm{o}}$ point by point for each argument on the horizontal axis, except possibly for points where $\mathrm{F}_{\mathrm{o}}$ jumps. (Recall that distribution functions are always continuous from the right, and except at jumps are continuous from the left. Since each jump contains a distinct rational number and the rationals are countable, there are at most a countable number of jumps. Then the set of jump points has Lebesgue measure zero, and there are continuity points arbitrarily close to any jump point. Because of right-continuity, distribution functions are uniquely determined by their values at their continuity points.) If $\mathbf{A}$ is an open set, then $\mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{Y}_{\mathrm{o}}$ implies $\liminf _{\mathrm{n} \rightarrow \infty} \mathrm{F}_{\mathrm{n}}(\mathbf{A}) \geq \mathrm{F}_{\mathrm{o}}(\mathbf{A})$; conversely, A closed implies limsup ${ }_{n \rightarrow \infty} \mathrm{~F}_{\mathrm{n}}(\mathbf{A}) \leq \mathrm{F}_{\mathrm{o}}(\mathbf{A})$ see P. Billingsley (1968), Theorem 2.1. Convergence in distribution is also called weak convergence in the space of distribution functions.
4.1.7. The relationships between different types of convergence are summarized in Figure 4.2. In this table, " $A=B$ " means that A implies $B$, but not vice versa, and "A $\Longleftrightarrow \Rightarrow B$ " means that $A$ and B are equivalent. Explanations and examples are given in Sections 4.1.8-4.1.18. On first reading, skim these sections and skip the proofs.
4.1.8. $\mathrm{Y}_{\mathrm{n}} \rightarrow_{\text {as }} \mathrm{Y}_{\mathrm{o}}$ implies $\operatorname{Prob}\left(\mathbf{W}_{\mathrm{n}}\right) \leq \operatorname{Prob}\left(\mathrm{U}_{m \geq n} \mathbf{W}_{\mathrm{m}}\right) \rightarrow 0$, and hence $\mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{p}} \mathrm{Y}_{\mathrm{o}}$. However,
$\operatorname{Prob}\left(\mathbf{W}_{\mathrm{n}}\right) \rightarrow 0$ does not necessarily imply that the probability of $U_{m \geq n} \mathbf{W}_{\mathrm{m}}$ is small, so $\mathrm{Y}_{\mathrm{n}} \vec{p}_{\mathrm{p}} \mathrm{Y}_{\mathrm{o}}$ does
not imply $Y_{n} \rightarrow_{\text {as }} Y_{0}$. For example, take the universe of states of nature to be the points on the unit circle with uniform probability, take the $\mathbf{W}_{\mathrm{n}}$ to be successive arcs of length $2 \pi / \mathrm{n}$, and take $\mathrm{Y}_{\mathrm{n}}$ to be 1 on $\mathbf{W}_{\mathrm{n}}, 0$ otherwise. Then $\mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{p}} 0$ since $\operatorname{Pr}\left(\mathrm{Y}_{\mathrm{n}} \neq 0\right)=1 / \mathrm{n}$, but $\mathrm{Y}_{\mathrm{n}}$ fails to converge almost surely to zero since the successive arcs wrap around the circle an infinite number of times, and every s in the circle is in an infinite number of $\mathbf{W}_{\mathrm{n}}$.
4.1.9. Suppose $Y_{n} \rightarrow_{p} Y_{0}$. It is a good exercise in manipulation of probabilities of events to show that $\mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{Y}_{\mathrm{o}}$. Given $\boldsymbol{\varepsilon}>0$, define $\mathbf{W}_{\mathrm{n}}$ as before to be the set of states of Nature where $\left|\mathrm{Y}_{\mathrm{n}}(\mathrm{s})-\mathrm{Y}_{\mathrm{o}}(\mathrm{s})\right|$ $>\varepsilon$. Given $y$, define $\mathbf{A}_{\mathrm{n}}, \mathbf{B}_{\mathrm{o}}$, and $\mathbf{C}_{\mathrm{o}}$ to be, respectively, the states of Nature with $\mathrm{Y}_{\mathrm{n}} \leq \mathrm{y}, \mathrm{Y}_{\mathrm{o}} \leq \mathrm{y}-\varepsilon$,
and $\mathrm{Y}_{\mathrm{o}} \leq \mathrm{y}+\varepsilon$. Then $\mathbf{B}_{\mathrm{o}} \subseteq \mathbf{A}_{\mathrm{n}} \cup \mathbf{W}_{\mathrm{n}}$ (i.e., $\mathrm{Y}_{\mathrm{o}}(\mathrm{s}) \leq \mathrm{y}-\boldsymbol{\varepsilon}$ implies either $\mathrm{Y}_{\mathrm{n}}(\mathrm{s}) \leq \mathrm{y}$ or $\left.\left|\mathrm{Y}_{\mathrm{o}}(\mathrm{s})-\mathrm{Y}_{\mathrm{n}}(\mathrm{s})\right|>\varepsilon\right)$ and $\mathbf{A}_{\mathrm{n}} \subseteq \mathbf{C}_{\mathrm{o}} \cup \mathbf{W}_{\mathrm{n}}$ (i.e., $\mathrm{Y}_{\mathrm{n}}(\mathrm{s}) \leq \mathrm{y}$ implies $\mathrm{Y}_{\mathrm{o}}(\mathrm{s}) \leq \mathrm{y}+\varepsilon$ or $\left|\mathrm{Y}_{\mathrm{o}}(\mathrm{s})-\mathrm{Y}_{\mathrm{n}}(\mathrm{s})\right|>\varepsilon$ ). Hence, for n large enough so $\operatorname{Prob}\left(\mathbf{W}_{\mathrm{n}}\right)<\varepsilon, \mathrm{F}_{\mathrm{o}}(\mathrm{y}-\boldsymbol{\varepsilon}) \equiv \operatorname{Prob}\left(\mathbf{B}_{\mathrm{o}}\right) \leq \operatorname{Prob}\left(\mathbf{A}_{\mathrm{n}}\right)+\operatorname{Prob}\left(\mathbf{W}_{\mathrm{n}}\right)<\mathrm{F}_{\mathrm{n}}(\mathrm{y})+\varepsilon$, and $\mathrm{F}_{\mathrm{n}}(\mathrm{y}) \equiv \operatorname{Prob}\left(\mathbf{A}_{\mathrm{n}}\right)$ $\leq \operatorname{Prob}\left(\mathbf{C}_{\mathrm{o}}\right)+\operatorname{Prob}\left(\mathbf{W}_{\mathrm{n}}\right)<\mathrm{F}_{\mathrm{o}}(\mathrm{Y}+\varepsilon)+\varepsilon$, implying $\mathrm{F}_{\mathrm{o}}(\mathrm{y}-\varepsilon)-\varepsilon \leq \lim _{\mathrm{n} \rightarrow \infty} \mathrm{F}_{\mathrm{n}}(\mathrm{y}) \leq \mathrm{F}_{\mathrm{o}}(\mathrm{y}+\varepsilon)+\varepsilon$. If y is a continuity point of $\mathrm{Y}_{\mathrm{o}}$, then $\mathrm{F}_{\mathrm{o}}(\mathrm{y}-\varepsilon)$ and $\mathrm{F}_{\mathrm{o}}(\mathrm{y}+\varepsilon)$ approach $\mathrm{F}_{\mathrm{o}}(\mathrm{y})$ as $\varepsilon \rightarrow 0$, implying $\lim _{\mathrm{n} \rightarrow \infty} \mathrm{F}_{\mathrm{n}}(\mathrm{y})=$ $F_{o}(y)$. This establishes that $Y_{n}{ }_{d} Y_{o}$.

Convergence in distribution of $Y_{n}$ to $Y_{o}$ does not imply that $Y_{n}$ and $Y_{o}$ are close to each other. For example, if $Y_{n}$ and $Y_{o}$ are i.i.d. standard normal, then $Y_{n}{ }_{d} Y_{o}$ trivially, but clearly not $Y_{n}{ }_{p} Y_{o}$ since $Y_{n}-Y_{o}$ is normal with variance 2 , and $\left|Y_{n}-Y_{o}\right|>\varepsilon$ with a positive, constant probability. However, there is a useful representation that is helpful in relating convergence in distribution and almost sure convergence; see P. Billingsley (1986), p. 343.

Theorem 4.1. (Skorokhod) If $Y_{n}{ }_{d} Y_{o}$, then there exist random variables $Y_{n}{ }^{\prime}$ and $Y_{o}{ }_{0}$ such that $\mathrm{Y}_{\mathrm{n}}$ and $\mathrm{Y}_{\mathrm{n}}{ }^{\prime}$ have the same CDF, as do $\mathrm{Y}_{\mathrm{o}}$ and $\mathrm{Y}_{\mathrm{o}}{ }^{\prime}$, and $\mathrm{Y}_{\mathrm{n}}{ }^{\prime} \overrightarrow{a s} \mathrm{Y}_{\mathrm{o}}{ }^{\prime}$.
4.1.10. Convergence in distribution and convergence in probability to a constant are equivalent. If $Y_{n} \rightarrow_{p} c$ constant, then $Y_{n}{ }_{d} c$ as a special case of 4.1.9 above. Conversely, $Y_{n}{ }_{d} c$ constant means $\mathrm{F}_{\mathrm{n}}(\mathrm{y}) \rightarrow \mathrm{F}_{\mathrm{p}}(\mathrm{y})$ at continuity points, where $\mathrm{F}_{\mathrm{c}}(\mathrm{y})=0$ for $\mathrm{y}<\mathrm{c}$ and $\mathrm{F}_{\mathrm{c}}(\mathrm{y})=1$ for $\mathrm{y} \geq \mathrm{c}$. Hence $\varepsilon>0$ implies $\operatorname{Prob}\left(\left|\mathrm{Y}_{\mathrm{n}}-\mathrm{c}\right|>\varepsilon\right)=\mathrm{F}_{\mathrm{n}}(\mathrm{c}-\varepsilon)+1-\mathrm{F}_{\mathrm{n}}(\mathrm{c}+\varepsilon) \rightarrow 0$, so $\mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{p}} \mathrm{c}$. This result implies particularly that the statements $Y_{n}-Y_{o} \rightarrow_{p} 0$ and $Y_{n}-Y_{o} \rightarrow_{d} 0$ are equivalent. Then, $Y_{n}-Y_{o} \rightarrow_{d} 0$ implies $Y_{n} \rightarrow_{d} Y_{o}$, but the reverse implication does not hold.
4.1.11. The condition that convergence in distribution is equivalent to convergence of expectations of all bounded continuous functions is a fundamental mathematical result called the Helly-Bray theorem. Intuitively, the reason the theorem holds is that bounded continuous functions can be approximated closely by sums of continuous "almost-step" functions, and the expectations of "almost step" functions closely approximate points of CDF's. A proof by J. Davidson (1994), p. 352, employs the Skorokhod representation theorem 4.1.
4.1.12. A Chebyshev-like inequality is obtained by noting for a random variable Z with density $\mathrm{f}(\mathrm{z})$ that $\mathbf{E}|\mathrm{Z}|^{\rho}=\int|\mathrm{z}|^{\rho} \mathrm{f}(\mathrm{z}) \mathrm{dz} \geq \int_{|z| \geq \varepsilon} \varepsilon^{\complement} \mathrm{f}(\mathrm{z}) \mathrm{dz}=\varepsilon^{\rho} \operatorname{Prob}(|\mathrm{Z}|>\varepsilon)$, or $\operatorname{Prob}(|\mathrm{Z}|>\varepsilon) \leq \mathbf{E}|\mathrm{Z}|^{\rho} / \varepsilon^{\rho}$.
(When $\rho=2$, this is the conventional Chebyshev inequality. When $\rho=1$, one has $\operatorname{Prob}(|Z|>\varepsilon) \leq$ $\mathbf{E}|\mathrm{Z}| / \varepsilon$.) Taking $\mathrm{Z}=\mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}}$, one has $\lim _{\mathrm{n}-\infty} \operatorname{Prob}\left(\left|\mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}}\right|>\varepsilon\right) \leq \boldsymbol{\varepsilon}^{-\rho} \cdot \lim _{\mathrm{n}-\infty} \mathbf{E}\left|\mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}}\right|^{\rho}$. Hence, convergence in $\rho$-mean (for any $\rho>0$ ) implies convergence in probability. However, convergence almost surely or in probability does not necessarily imply convergence in $\rho$-mean. Suppose the sample space is the unit interval with uniform probability, and $Y_{n}(s)=e^{n^{\prime}}$ for $\mathrm{s} \leq \mathrm{n}^{-2}$, zero otherwise. Then $\mathrm{Y}_{\mathrm{n}} \overrightarrow{a s} 0$ since $\operatorname{Prob}\left(\mathrm{Y}_{\mathrm{m}} \neq 0\right.$ for any $\left.\mathrm{m}>\mathrm{n}\right) \leq \mathrm{n}^{-2}$, but $\mathbf{E}\left|\mathrm{Y}_{\mathrm{n}}\right|^{\rho}=\mathrm{e}^{\rho \mathrm{n}} / \mathrm{n}^{2} \rightarrow+\infty$ for any $\rho>0$.

## FIGURE 4.2. RELATIONS BETWEEN STOCHASTIC LIMITS

(Section numbers for details are given in parentheses)

| 1 | $\mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{as}} \mathrm{Y}_{\mathrm{o}} \quad \stackrel{(1.8)}{\rightleftharpoons} \quad \mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{p}} \mathrm{Y}_{\mathrm{o}} \quad \stackrel{(1.9)}{\rightleftharpoons} \quad \mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{Y}_{\mathrm{o}}$ |
| :---: | :---: |
|  | $\Uparrow$ <br> 介 <br> $\Uparrow$ <br> (1.4) <br> (1.3) <br> (1.10) <br> $\Downarrow$ <br> $\Downarrow$ <br> II |
| 2 | $\mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}} \rightarrow_{\mathrm{as}} 0 \underset{(1.8)}{\underset{ }{\Longrightarrow}} \mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}} \rightarrow_{\mathrm{p}} 0 \underset{(1.10)}{\Longleftrightarrow} \underset{\mathrm{n}}{\Longrightarrow}-\mathrm{Y}_{\mathrm{o}} \rightarrow_{\mathrm{d}} 0$ |
| 3 | $\mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{c}\left(\mathrm{a}\right.$ constant) $\Longleftrightarrow \Longleftrightarrow \mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{p}} \mathrm{c}$ |
| 4 | $\mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{Y}_{\mathrm{o}} \Longleftrightarrow \Longleftrightarrow \mathbf{E g}\left(\mathrm{Y}_{\mathrm{n}}\right) \rightarrow \mathbf{E g}\left(\mathrm{Y}_{\mathrm{o}}\right)$ for all bounded continuous g |
| 5 | $\left\\|\mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}}\right\\|_{\mathrm{\rho}} \rightarrow 0$ for some $\rho>0 \Longrightarrow \mathrm{Y}_{\mathrm{n}} \rightarrow \mathrm{p}_{\mathrm{p}} \mathrm{Y}_{\mathrm{o}} \quad$ (1.12) |
| 6 | $\left\\|\mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}}\right\\|_{\mathrm{p}} \leq \mathrm{M}($ all n$) \& \mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{p}} \mathrm{Y}_{\mathrm{o}} \xlongequal{\Longrightarrow}\left\\|\mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}}\right\\|_{\lambda} \rightarrow 0$ for $0<\lambda<\rho$ |
| 7 | $\mathrm{Y}_{\mathrm{n}} \rightarrow \mathrm{p}_{\mathrm{p}} \mathrm{Y}_{\mathrm{o}} \mid \Longrightarrow Y_{n_{k}} \rightarrow_{\text {as }} \mathrm{Y}_{\mathrm{o}}$ for some subsequence $\mathrm{n}_{\mathrm{k}}, \mathrm{k}=1,2, \ldots$ |
| 8 | $\sum_{n=1}^{\infty} \mathrm{P}\left(\left\|\mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}}\right\|>\varepsilon\right)<+\infty$ for each $\varepsilon>0 \Longrightarrow \mathrm{Y}_{\mathrm{n}} \rightarrow_{\text {as }} \mathrm{Y}_{\mathrm{o}}$ |
| 9 | $\sum_{n=1}^{\infty} \quad \mathbf{E}\left\|\mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}}\right\|^{\rho}<+\infty($ for some $\rho>0) \Longrightarrow \Longrightarrow \mathrm{Y}_{\mathrm{n}} \rightarrow_{\text {as }} \mathrm{Y}_{\text {o }} \quad$ (1.15) |
| 10 | $\mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{Y}_{\mathrm{o}} \& \mathrm{Z}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{p}} 0 \models \mathrm{Z}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{Y}_{\mathrm{o}} \quad$ (1.16) |
| 11 | $\mathrm{Y}_{\mathrm{n}} \vec{p} \mathrm{Y}_{\mathrm{o}} \\| \Rightarrow \mathrm{g}\left(\mathrm{Y}_{\mathrm{n}}\right) \vec{p}^{\mathrm{g}} \mathrm{g}\left(\mathrm{Y}_{\mathrm{o}}\right)$ for all continuous g (1.17) |
| 12 | $Y_{n} \rightarrow_{d} Y_{o} \xlongequal{\prime} \Rightarrow g\left(Y_{n}\right) \rightarrow_{d} g\left(Y_{o}\right)$ for all continuous $g$ (1.18) |

4.1.13. Adding a condition of a uniformly bounded $\rho$-order mean $\mathbf{E}\left|Y_{n}\right|^{\rho} \leq M$ to convergence in probability $\mathrm{Y}_{\mathrm{n}} \overrightarrow{\mathrm{p}} \mathrm{Y}_{\mathrm{o}}$ yields the result that $\mathbf{E}\left|\mathrm{Y}_{\mathrm{o}}\right|^{\lambda}$ exists for $0<\lambda \leq \rho$, and $\mathbf{E}\left|\mathrm{Y}_{\mathrm{n}}\right|^{\lambda} \rightarrow \mathbf{E}\left|\mathrm{Y}_{\mathrm{o}}\right|^{\lambda}$ for 0 $<\lambda<\rho$. This result can be restated as "the moments of the limit equal the limit of the moments" for moments of order $\lambda$ less than $\rho$. Replacing $Y_{n}$ by $Y_{n}-Y_{o}$ and $Y_{o}$ by 0 gives the result in Figure 4.2.

To prove these results, we will find useful the property of moments that $\mathbf{E}|\mathrm{Y}|^{\lambda} \leq\left(\mathbf{E}|\mathrm{Y}|^{\rho}\right)^{\lambda / \rho}$ for $0<\lambda<\rho$. (This follows from Holder's inequality (2.1.11), which states $\mathbf{E}|\mathrm{UV}| \leq\left(\mathbf{E}|\mathrm{U}|^{\mathrm{r}}\right)^{1 / \mathrm{r}}\left(\mathbf{E}|\mathrm{V}|^{\mathrm{s}}\right)^{1 / \mathrm{s}}$ for $\mathrm{r}, \mathrm{s}>0$ and $\mathrm{r}^{-1}+\mathrm{s}^{-1}=1$, by taking $\mathrm{U}=|\mathrm{Y}|^{\lambda}, \mathrm{V}=1$, and $\mathrm{r}=\rho / \lambda$.) An immediate implication is $\mathbf{E}\left|\mathrm{Y}_{\mathrm{n}}\right|^{\lambda} \leq \mathrm{M}^{\lambda / \rho}$. Define $\mathrm{g}(\mathrm{y}, \lambda, \mathrm{k})=\min \left(|\mathrm{y}|^{\lambda}, \mathrm{k}^{\lambda}\right)$, and note that since it is continuous and bounded, the Healy-Bray theorem implies $\operatorname{Eg}\left(\mathrm{Y}_{\mathrm{n}}, \lambda, \mathrm{k}\right) \rightarrow \mathbf{E g}\left(\mathrm{Y}_{\mathrm{o}}, \lambda, \mathrm{k}\right)$. Therefore,

$$
\begin{aligned}
\mathrm{M}^{\lambda_{\rho}} \geq \mathbf{E}\left|\mathrm{Y}_{\mathrm{n}}\right|^{\lambda} \geq \mathbf{E g}\left(\mathrm{Y}_{\mathrm{n}}, \lambda, \mathrm{k}\right)= & \int_{-k}^{k}|\mathrm{y}|^{\lambda} \mathrm{f}_{\mathrm{n}}(\mathrm{y}) \mathrm{dy}+\mathrm{k}^{\lambda} \cdot \operatorname{Prob}\left(\left|\mathrm{Y}_{\mathrm{n}}\right|>\mathrm{k}\right) \\
& \rightarrow \int_{-k}^{k}|\mathrm{y}|^{\lambda} \mathrm{f}_{\mathrm{o}}(\mathrm{y}) \mathrm{dy}+\mathrm{k}^{\lambda} \operatorname{Prob}\left(\left|\mathrm{Y}_{\mathrm{o}}\right|>\mathrm{k}\right) .
\end{aligned}
$$

Letting $\mathrm{k} \rightarrow \infty$ establishes that $\mathbf{E}\left|\mathrm{Y}_{\mathrm{o}}\right|^{\lambda}$ exists for $0<\lambda \leq \rho$. Further, for $\lambda<\rho$,

$$
0 \leq \mathbf{E}\left|\mathrm{Y}_{\mathrm{n}}\right|^{\lambda}-\mathbf{E g}\left(\mathrm{Y}_{\mathrm{n}}, \lambda, \mathrm{k}\right) \leq \int_{|y|>k}|\mathrm{y}|^{\lambda} \mathrm{f}_{\mathrm{n}}(\mathrm{y}) \mathrm{dy} \leq \mathrm{k}^{\lambda-\rho} \int_{|y|>k}|\mathrm{y}|^{\rho} \mathrm{f}_{\mathrm{n}}(\mathrm{y}) \mathrm{dy} \leq \mathrm{k}^{\lambda-\rho} \mathrm{M} .
$$

Choose k sufficiently large so that $\mathrm{k}^{\lambda-\rho} \mathrm{M}<\varepsilon$. The same inequality holds for $\mathrm{Y}_{\mathrm{o}}$. Choose n sufficiently large so that $\left|\mathbf{E g}\left(\mathrm{Y}_{\mathrm{n}}, \lambda, \mathrm{k}\right) \rightarrow \mathbf{E g}\left(\mathrm{Y}_{\mathrm{o}}, \lambda, \mathrm{k}\right)\right|<\varepsilon$. Then

$$
\left.|\mathbf{E}| \mathrm{Y}_{\mathrm{n}}\right|^{\lambda}-\mathbf{E}\left|\mathrm{Y}_{\mathrm{o}}\right|^{\lambda}\left|\leq|\mathbf{E}| \mathrm{Y}_{\mathrm{n}}\right|^{\lambda}-\mathbf{E g}\left(\mathrm{Y}_{\mathrm{n}}\right)\left|+\left|\mathbf{E g}\left(\mathrm{Y}_{\mathrm{n}}\right)-\mathbf{E g}\left(\mathrm{Y}_{\mathrm{o}}\right)\right|+\left|\mathbf{E g}\left(\mathrm{Y}_{\mathrm{o}}\right)-\mathbf{E}\right| \mathrm{Y}_{\mathrm{o}}\right|^{\lambda} \mid \leq 3 \varepsilon .
$$

This proves that $\mathbf{E}\left|\mathrm{Y}_{\mathrm{n}}\right|^{\lambda} \rightarrow \mathbf{E}\left|\mathrm{Y}_{\mathrm{o}}\right|^{\lambda}$.
An example shows that $\mathbf{E}\left|Z_{n}\right|^{\lambda} \rightarrow 0$ for $\lambda<\rho$ does not imply $\mathbf{E}\left|Z_{n}\right|^{\rho}$ bounded. Take $Z_{n}$ discrete with support $\{0, \mathrm{n}\}$ and probability $\log (\mathrm{n}) / \mathrm{n}$ at n . Then for $\lambda<1, \mathbf{E}\left|\mathrm{Z}_{\mathrm{n}}\right|^{\lambda}=\log (\mathrm{n}) / \mathrm{n}^{1-\lambda} \rightarrow 0$, but $\mathbf{E}\left|\mathrm{Z}_{\mathrm{n}}\right|^{1}$ $=\log (\mathrm{n}) \rightarrow+\infty$.
4.1.14. If $\mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{p}} \mathrm{Y}_{\mathrm{o}}$, then $\operatorname{Prob}\left(\mathbf{W}_{\mathrm{n}}\right) \rightarrow 0$. Choose a subsequence $\mathrm{n}_{\mathrm{k}}$ such that $\operatorname{Prob}\left(\boldsymbol{W}_{n_{k}}\right) \leq 2^{-\mathrm{k}}$. Then $\operatorname{Prob}\left(\bigcup_{k^{\prime}>k} \boldsymbol{W}_{n_{k}}\right) \leq \sum_{k^{\prime}>k} \operatorname{Prob}\left(\boldsymbol{W}_{n_{k}}\right) \leq \sum_{k^{\prime}>k} 2^{-\mathrm{k}^{\prime}}=2^{-\mathrm{k}}$, implying $Y_{n_{k}} \rightarrow_{\mathrm{as}} \mathrm{Y}_{\mathrm{o}}$.
4.1.15. Conditions for a.s. convergence follow from this basic probability theorem:

Theorem 4.2. (Borel-Cantelli) If $\mathbf{A}_{\mathrm{i}}$ is any sequence of events in a probability space ( $\mathbf{S}, \boldsymbol{F}, \mathrm{P}$ ), $\sum_{n=1}^{\infty} \mathrm{P}\left(\mathbf{A}_{\mathrm{i}}\right)<+\infty$ implies that almost surely only a finite number of the events $\mathbf{A}_{\mathrm{i}}$ occur. If $\mathbf{A}_{\mathrm{i}}$ is a sequence of independent events, then $\sum_{n=1}^{\infty} \quad \mathrm{P}\left(\mathbf{A}_{\mathrm{i}}\right)=+\infty$ implies that almost surely an infinite number of the events $\mathbf{A}_{\mathrm{i}}$ occur.

Apply the Borel-Cantelli theorem to the events $\mathbf{A}_{\mathrm{i}}=\left\{\mathrm{s} \in \mathbf{S}| | \mathrm{Y}_{\mathrm{i}}-\mathrm{Y}_{\mathrm{o}} \mid>\varepsilon\right\}$ to conclude that $\sum_{n=1}^{\infty} \mathrm{P}\left(\mathbf{A}_{\mathrm{i}}\right)<+\infty$ implies that almost surely only a finite number of the events $\mathbf{A}_{\mathrm{i}}$ occur, and hence $\left|Y_{i}-Y_{o}\right| \leq \varepsilon$ for all i sufficiently large. Thus, $Y_{n}-Y_{o} \rightarrow_{\text {as }} 0$, or $Y_{n} \rightarrow_{a s} Y_{o}$. For the next result in the table, use (1.12) to get $\operatorname{Prob}\left(\bigcup_{m \geq n} \boldsymbol{W}_{m}\right) \leq \sum_{m>n} \operatorname{Prob}\left(\mathbf{W}_{\mathrm{m}}\right) \leq \boldsymbol{\varepsilon}^{-\rho} \sum_{m>n} \mathbf{E}\left|\mathrm{Y}_{\mathrm{m}}-\mathrm{Y}_{\mathrm{o}}\right|^{\rho}$. Apply Theorem 4.2 to conclude that if this right-hand expression is finite, then $Y_{n} \rightarrow_{\text {as }} Y_{o}$. The example at the end of (1.12) shows that almost sure convergence does not imply convergence in $\rho$-mean. Also, the example mentioned in 1.8 which has convergence in probability but not almost sure convergence can be constructed to have $\rho$-mean convergence but not almost sure convergence.
4.1.16. A result termed the Slutsky theorem which is very useful in applied work is that if two random variables $Y_{n}$ and $Z_{n}$ have a difference which converges in probability to zero, and if $Y_{n}$ converges in distribution to $\mathrm{Y}_{\mathrm{o}}$, then $\mathrm{Z}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{Y}_{\mathrm{o}}$ also. In this case, $\mathrm{Y}_{\mathrm{n}}$ and $\mathrm{Z}_{\mathrm{n}}$ are termed asymptotically equivalent. The argument demonstrating this result is similar to that for 4.1.9. Let $\mathrm{F}_{\mathrm{n}}$ and $\mathrm{G}_{\mathrm{n}}$ be the CDF's of $\mathrm{Y}_{\mathrm{n}}$ and $\mathrm{Z}_{\mathrm{n}}$ respectively. Let y be a continuity point of $\mathrm{F}_{\mathrm{o}}$ and define the following events:

$$
\mathbf{A}_{\mathrm{n}}=\left\{\mathrm{s} \mid \mathrm{Z}_{\mathrm{n}}(\mathrm{~s})<\mathrm{y}\right\}, \mathbf{B}_{\mathrm{n}}=\left\{\mathrm{s} \mid \mathrm{Y}_{\mathrm{n}}(\mathrm{~s}) \leq \mathrm{y}-\varepsilon\right\}, \mathbf{C}_{\mathrm{n}}=\left\{\mathrm{s} \mid \mathrm{Y}_{\mathrm{n}}(\mathrm{~s}) \leq \mathrm{y}+\varepsilon\right\}, \mathbf{D}_{\mathrm{n}}=\left\{\mathrm{s}| | \mathrm{Y}_{\mathrm{n}}(\mathrm{~s})-\mathrm{Z}_{\mathrm{n}}(\mathrm{~s}) \mid>\varepsilon\right\} .
$$

Then $\mathbf{A}_{\mathrm{n}} \subseteq \mathbf{C}_{\mathrm{n}} \cup \mathbf{D}_{\mathrm{n}}$ and $\mathbf{B}_{\mathrm{n}} \subseteq \mathbf{A}_{\mathrm{n}} \cup \mathbf{D}_{\mathrm{n}}$, implying $\mathrm{F}_{\mathrm{n}}(\mathrm{y}-\varepsilon)-\operatorname{Prob}\left(\mathbf{D}_{\mathrm{n}}\right) \leq \mathrm{G}_{\mathrm{n}}(\mathrm{y}) \leq \mathrm{F}_{\mathrm{n}}(\mathrm{y}+\varepsilon)+\operatorname{Prob}\left(\mathbf{D}_{\mathrm{n}}\right)$. Given $\delta>0$, one can choose $\varepsilon>0$ such that $\mathrm{y}-\varepsilon$ and $\mathrm{y}+\varepsilon$ are continuity points of $\mathrm{F}_{\mathrm{n}}$, and such that $\mathrm{F}_{\mathrm{o}}(\mathrm{y}+\varepsilon)-\mathrm{F}_{\mathrm{o}}(\mathrm{y}-\varepsilon)<\delta / 3$. Then one can choose n sufficiently large so that $\operatorname{Prob}\left(\mathbf{D}_{\mathrm{n}}\right)<\delta / 3, \mid \mathrm{F}_{\mathrm{n}}(\mathrm{y}+\varepsilon)-$ $\mathrm{F}_{\mathrm{o}}(\mathrm{y}+\varepsilon) \mid<\delta / 3$ and $\left|\mathrm{F}_{\mathrm{n}}(\mathrm{y}+\varepsilon)-\mathrm{F}_{\mathrm{o}}(\mathrm{y}+\varepsilon)\right|<\delta / 3$. Then $\left|\mathrm{G}_{\mathrm{n}}(\mathrm{y})-\mathrm{F}_{\mathrm{o}}(\mathrm{y})\right|<\delta$.
4.1.17 A useful property of convergence in probability is the following result:

Theorem 4.3. (Continuous Mapping Theorem) If $g(y)$ is a continuous function on an open set containing the support of $Y_{o}$, then $Y_{n} \rightarrow_{p} Y_{o}$ implies $g\left(Y_{n}\right) \rightarrow_{p} g\left(Y_{o}\right)$. The result also holds for vectors of random variables, and specializes to the rules that if $\mathrm{Y}_{1 \mathrm{n}} \overrightarrow{\mathrm{p}}_{\mathrm{p}} \mathrm{Y}_{10}, \mathrm{Y}_{2 \mathrm{n}} \rightarrow_{\mathrm{p}} \mathrm{Y}_{20}$, and $\mathrm{Y}_{3 \mathrm{n}} \rightarrow_{\mathrm{p}} \mathrm{Y}_{30}$ then (a) $\mathrm{Y}_{1 \mathrm{n}} \cdot \mathrm{Y}_{2 \mathrm{n}}+\mathrm{Y}_{3 \mathrm{n}} \overrightarrow{\mathrm{p}} \mathrm{Y}_{10} \cdot \mathrm{Y}_{20}+\mathrm{Y}_{30}$, and (b) if $\operatorname{Prob}\left(\left|\mathrm{Y}_{20}\right|<\varepsilon\right)=0$ for some $\varepsilon>0$, then $\mathrm{Y}_{1 \mathrm{n}} / \mathrm{Y}_{2 \mathrm{n}} \overrightarrow{\mathrm{p}}$ $\mathrm{Y}_{10} / \mathrm{Y}_{20}$. In these limits, $\mathrm{Y}_{10}, \mathrm{Y}_{20}$, and/or $\mathrm{Y}_{30}$ may be constants.

Proof: Given $\varepsilon>0$, choose M such that $\mathrm{P}\left(\left|\mathrm{Y}_{\mathrm{o}}\right|>\mathrm{M}\right)<\varepsilon$. Let $\mathbf{A}_{\mathrm{o}}$ be the set of y in the support of $\mathrm{Y}_{\mathrm{o}}$ that satisfy $|\mathrm{y}| \leq \mathrm{M}$. Then $\mathbf{A}_{\mathrm{o}}$ is compact. Mathematical analysis can be used to show that there exists a nested sequence of sets $\mathbf{A}_{0} \subseteq \mathbf{A}_{1} \subseteq \mathbf{A}_{2} \subseteq \mathbf{A}_{3}$ with $\mathbf{A}_{3}$ an open neighborhood of $\mathbf{A}_{o}$ on which g is continuous, $\mathbf{A}_{2}$ compact, and $\mathbf{A}_{1}$ open. From 4.16, $\liminf _{\mathrm{n} \rightarrow \infty} \mathrm{F}_{\mathrm{n}}\left(\mathbf{A}_{1}\right) \geq \mathrm{F}_{\mathrm{o}}\left(\mathbf{A}_{1}\right) \geq 1-\varepsilon$ implies there exists $\mathrm{n}_{1}$ such that for $\mathrm{m}>\mathrm{n}_{1}, \mathrm{~F}_{\mathrm{m}}\left(\mathbf{A}_{1}\right) \geq 1-2 \varepsilon$. The continuity of g implies that for each $\mathrm{y} \in \mathbf{A}_{2}$, there exists $\delta_{\mathrm{y}}>0$ such that $\left|\mathrm{y}^{\prime}-\mathrm{y}\right|<\delta_{\mathrm{y}} \Rightarrow\left|\mathrm{g}\left(\mathrm{y}^{\prime}\right)-\mathrm{g}(\mathrm{y})\right|<\varepsilon$. These $\delta_{\mathrm{y}}$-neighborhoods cover $\mathbf{A}_{2}$. Then $\mathbf{A}_{2}$ has a finite subcover. Let $\delta$ be the smallest value of $\delta_{\mathrm{y}}$ in this finite subcover. Then, g is uniformly continuous: $\mathrm{y} \in \mathbf{A}_{2}$ and $\left|\mathrm{y}^{\prime}-\mathrm{y}\right|<\delta$ imply $\left|\mathrm{g}\left(\mathrm{y}^{\prime}\right)-\mathrm{g}(\mathrm{y})\right|<\varepsilon$. Choose $\mathrm{n}>\mathrm{n}_{1}$ such that for $\mathrm{m}>\mathrm{n}, \mathrm{P}\left(\mid \mathrm{Y}_{\mathrm{m}}\right.$ $\left.-\mathrm{Y}_{\mathrm{o}} \mid>\delta\right)<\varepsilon / 2$. Then for $\mathrm{m}>\mathrm{n}, \mathrm{P}\left(\left|\mathrm{g}\left(\mathrm{Y}_{\mathrm{m}}\right)-\mathrm{g}\left(\mathrm{Y}_{\mathrm{o}}\right)\right|>\varepsilon\right) \leq \mathrm{P}\left(\left|\mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}}\right|>\delta\right)+\mathrm{P}\left(\left|\mathrm{Y}_{\mathrm{o}}\right|>\mathrm{M}\right)+1-$ $\mathrm{F}_{\mathrm{m}}\left(\mathrm{A}_{1}\right) \leq 4 \boldsymbol{\varepsilon}$.
4.1.18 The preceding result has an analog for convergence in distribution. This result establishes, for example, that if $Y_{n} \rightarrow_{d} Y_{o}$, with $Y_{o}$ standard normal and $g(y)=y^{2}$, then $Y_{o}$ is chi-squared, so that that $Y_{n}{ }^{2}$ converges in distribution to a chi-squared random variable.

Theorem 4.4. If $g(y)$ is a continuous function on an open set containing the support of $Y_{o}$, then $Y_{n} \rightarrow_{d} Y_{o}$ implies $g\left(Y_{n}\right) \rightarrow_{d} g\left(Y_{o}\right)$. The result also holds for vectors of random variables.

Proof: The Skorokhod representation given in Theorem 4.1 implies there exist $\mathrm{Y}_{\mathrm{n}}{ }^{\prime}$ and $\mathrm{Y}_{\mathrm{o}}{ }^{\prime}$ that have the same distributions as $Y_{n}$ and $Y_{0}$, respectively, and satisfy $Y_{n}{ }^{\prime}{ }_{\text {as }} Y_{o}{ }^{\prime}$. Then, Theorem 4.3 implies $\mathrm{g}\left(\mathrm{Y}_{\mathrm{n}}{ }^{\prime}\right) \rightarrow_{\text {as }} \mathrm{g}\left(\mathrm{Y}_{\mathrm{o}}{ }^{\prime}\right)$, and results 4.1.8 and 4.1.9 above then imply $\mathrm{g}\left(\mathrm{Y}_{\mathrm{n}}{ }^{\prime}\right) \rightarrow_{\mathrm{d}} \mathrm{g}\left(\mathrm{Y}_{\mathrm{o}}{ }^{\prime}\right)$. Because of the common distributions, this is the result in Theorem 4.4. For this reason, this result is also sometimes referred to as (part of) the continuous mapping theorem. The Slutsky theorem, result 4.1.10, is a special case of the continous mapping Theorems 4.3 and 4.4. For clarity, I also give a direct proof of Theorem 4.4. Construct the sets $\mathbf{A}_{\mathrm{o}} \subseteq \mathbf{A}_{1} \subseteq \mathbf{A}_{2} \subseteq \mathbf{A}_{3}$ as in the proof of Theorem 4.3. A theorem from mathematical analysis (Urysohn) states that there exists a continuous function $r$ with values between zero and one that satisfies $r(y)=1$ for $y \in \mathbf{A}_{1}$ and $r(y)=0$ for $y \notin \mathbf{A}_{3}$. Then $g^{*}(y)=g(y) \cdot r(y)$ is continuous everywhere. From the Healy-Bray theorem, $\mathrm{Y}_{\mathrm{n}}{ }_{\mathrm{d}} \mathrm{Y}_{\mathrm{o}} \Longleftrightarrow \boldsymbol{E} h\left(\mathrm{Y}_{\mathrm{n}}\right) \rightarrow \mathbf{E} h\left(\mathrm{Y}_{\mathrm{o}}\right)$ for all continuous bounded $\mathrm{h} \Longrightarrow \mathbf{E h}\left(\mathrm{g}^{*}\left(\mathrm{Y}_{\mathrm{n}}\right)\right) \rightarrow \mathbf{E h}\left(\mathrm{g}^{*}\left(\mathrm{Y}_{\mathrm{o}}\right)\right)$ for all continuous bounded h , since the composition of continuous bounded functions is continuous and bounded $\Longleftrightarrow=\Rightarrow g^{*}\left(Y_{n}\right) \rightarrow_{d} g^{*}\left(Y_{o}\right)$.

But $\mathrm{P}\left(\mathrm{g}^{*}\left(\mathrm{Y}_{\mathrm{n}}\right) \neq \mathrm{g}\left(\mathrm{Y}_{\mathrm{n}}\right)\right) \leq \mathrm{P}\left(\mathrm{Y}_{\mathrm{n}} \notin \mathbf{A}_{1}\right) \leq 2 \varepsilon$ for n sufficiently large, and $\mathrm{g}^{*}\left(\mathrm{Y}_{\mathrm{o}}\right)=\mathrm{g}\left(\mathrm{Y}_{\mathrm{o}}\right)$. Then, 4.1.16 and $g^{*}\left(Y_{n}\right)-g\left(Y_{n}\right) \rightarrow_{p} 0$ imply $g^{*}\left(Y_{n}\right) \rightarrow_{d} g^{*}\left(Y_{o}\right)$.
4.1.19. Convergence properties are sometimes summarized in a notation called $O_{\mathrm{p}}(\cdot)$ and $o_{\mathrm{p}}(\cdot)$ which is very convenient for manipulation. (Sometimes too convenient; it is easy to get careless and make mistakes using this calculus.) The definition of $o_{\mathrm{p}}(\cdot)$ is that a random sequence $\mathrm{Y}_{\mathrm{n}}$ is $o_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right)$ if $\mathrm{n}^{-\alpha} \mathrm{Y}_{\mathrm{n}}$ converges in probability to zero; and one then writes $\mathrm{Y}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right)$. Then, $\mathrm{Y}_{\mathrm{n}}{ }_{\mathrm{p}} \mathrm{Y}_{\mathrm{o}}$ is also written $\mathrm{Y}_{\mathrm{n}}=\mathrm{Y}_{\mathrm{o}}+o_{\mathrm{p}}(1)$, and more generally $\mathrm{n}^{-\alpha}\left(\mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}}\right) \rightarrow_{\mathrm{p}} 0$ is written $\mathrm{Y}_{\mathrm{n}}-\mathrm{Y}_{\mathrm{o}}=o_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right)$. Thus $o_{\mathrm{p}}(\cdot)$ is a notation for convergence in probability to zero of a suitably normalized sequence of random variables. When two sequences of random variables $Y_{n}$ and $Z_{n}$ are asymptotically equivalent, so that they satisfy $\mathrm{Y}_{\mathrm{n}}-\mathrm{Z}_{\mathrm{n}}=o_{\mathrm{p}}(1)$, then they have a common limiting distribution by Slutsky's theorem, and this is sometime denoted $\mathrm{Y}_{\mathrm{n}} \sim_{\mathrm{a}} \mathrm{Z}_{\mathrm{n}}$.

The notation $\mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}(1)$ is defined to mean that given $\varepsilon>0$, there exists a large M (not depending on $n$ ) such that $\operatorname{Prob}\left(\left|\mathrm{Y}_{\mathrm{n}}\right|>\mathrm{M}\right)<\varepsilon$ for all n . A sequence with this property is called stochastically bounded. More generally, $\mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right)$ means $\operatorname{Prob}\left(\left|\mathrm{Y}_{\mathrm{n}}\right|>\mathrm{M} \cdot \mathrm{n}^{\alpha}\right)<\varepsilon$ for all n . A sequence that is convergent in distribution is stochastically bounded: If $Y_{n} \rightarrow_{d} Y_{o}$, then one can find $M$ and $n_{o}$ such that $\pm M$ are continuity points of $Y_{o}, \operatorname{Prob}\left(\left|Y_{o}\right| \leq M\right)>1-\varepsilon / 2,\left|F_{n}(M)-F_{0}(M)\right|<\varepsilon / 4$ and $\mid F_{n}(-M)-F_{0}(-$ $\mathrm{M}) \mid<\varepsilon / 4$ for $\mathrm{n}>\mathrm{n}_{\mathrm{o}}$. Then $\operatorname{Prob}\left(\left|\mathrm{Y}_{\mathrm{n}}\right|>\mathrm{M}\right)<\varepsilon$ for $\mathrm{n}>\mathrm{n}_{\mathrm{o}}$. This implies $\mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}(1)$. On the other hand, one can have $\mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}(1)$ without having convergence to any distribution (e.g., consider $\mathrm{Y}_{\mathrm{n}} \equiv 0$ for n odd and $\mathrm{Y}_{\mathrm{n}}$ standard normal for n even). The notation $\mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right)$ means $\mathrm{n}^{-\alpha} \mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}(1)$.

Most of the properties of $O_{\mathrm{p}}(\cdot)$ and $\sigma_{\mathrm{p}}(\cdot)$ are obvious restatements of results from Figure 4.2. For example, $\mathrm{n}^{-\alpha} \mathrm{Y}_{\mathrm{n}}=o_{\mathrm{p}}(1)$, or $\mathrm{n}^{-\alpha} \mathrm{Y}_{\mathrm{n}} \rightarrow_{\mathrm{p}} 0$, immediately implies for any $\varepsilon>0$ that there exists $\mathrm{n}_{\mathrm{o}}$ such that for $\mathrm{n}>\mathrm{n}_{0}, \operatorname{Prob}\left(\left|\mathrm{n}^{-0} \mathrm{Y}_{\mathrm{n}}\right|>\varepsilon\right)<\varepsilon$. For each $\mathrm{n} \leq \mathrm{n}_{0}$, one can find $\mathrm{M}_{\mathrm{n}}$ such that $\operatorname{Prob}\left(\left|\mathrm{n}^{-0} \mathrm{Y}_{\mathrm{n}}\right|>\mathrm{M}_{\mathrm{n}}\right)<\varepsilon$. Then, taking $M$ to be the maximum of $\varepsilon$ and the $M_{n}$ for $n \leq n_{o}$, one has $\operatorname{Prob}\left(\left|n^{-\alpha} Y_{n}\right|>M\right)<\varepsilon$ for all n , and hence $\mathrm{n}^{-\alpha} \mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}(1)$. The results above can be summarized in the following string of inplications:


An abbreviated list of rules for $o_{\mathrm{p}}$ and $O_{\mathrm{p}}$ is given in Figure 4.3. We prove the very useful rule 6 in this figure: Given $\varepsilon>0, \mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right) \Longrightarrow \exists \mathrm{M}>0$ such that $\operatorname{Prob}\left(\left|\mathrm{n}^{-\alpha} \mathrm{Y}_{\mathrm{n}}\right|>\mathrm{M}\right)<\varepsilon / 2$. Next $\mathrm{Z}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\beta}\right)$ implies $\exists \mathrm{n}_{\mathrm{o}}$ such that for $\mathrm{n}>\mathrm{n}_{\mathrm{o}}, \operatorname{Prob}\left(\left|\mathrm{n}^{-\beta} \mathrm{Z}_{\mathrm{n}}\right|>\varepsilon / \mathrm{M}\right)<\varepsilon / 2$. Hence $\operatorname{Prob}\left(\left|\mathrm{n}^{-\alpha \cdot \beta} \mathrm{Y}_{\mathrm{n}} \mathrm{Z}_{\mathrm{n}}\right|>\right.$ $\varepsilon) \leq \operatorname{Prob}\left(\left|\mathrm{n}^{-\alpha} \mathrm{Y}_{\mathrm{n}}\right|>\mathrm{M}\right)+\operatorname{Prob}\left(\left|\mathrm{n}^{-\beta} \mathrm{Z}_{\mathrm{n}}\right|>\varepsilon / \mathrm{M}\right)<\varepsilon$. Demonstration of the remaining rules is left as an exercise.

## FIGURE 4.3. RULES FOR $O_{\mathrm{p}}(\cdot)$ AND $o_{\mathrm{p}}(\cdot)$

|  | Definition: $\mathbf{Y}_{\mathbf{n}}=o_{\mathrm{p}}\left(\mathbf{n}^{\alpha}\right) \Longrightarrow \operatorname{Prob}\left(\left\|\mathbf{n}^{-\alpha} \mathbf{Y}_{\mathrm{n}}\right\|>\varepsilon\right) \rightarrow 0$ for each $\varepsilon>0$. <br> Definition: $\mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right) \Longrightarrow$ for each $\varepsilon>0$, there exists $\mathrm{M}>0$ such that $\operatorname{Prob}\left(\left\|n^{-\alpha} Y_{n}\right\|>M\right)<\varepsilon$ for all $n$ |
| :---: | :---: |
| 1 | $\mathrm{Y}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right) \Longrightarrow \Rightarrow \mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right)$ |
| 2 | $\mathrm{Y}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right) \& \beta>\alpha \Longrightarrow \mathrm{Y}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\beta}\right)$ |
| 3 | $\mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right) \& \beta>\alpha=\mathrm{Y}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\beta}\right)$ |
| 4 | $\mathrm{Y}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right) \& \mathrm{Z}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\beta}\right) \Longrightarrow \mathrm{Y}_{\mathrm{n}} \cdot \mathrm{Z}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\alpha+\beta}\right)$ |
| 5 | $\mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right) \& \mathrm{Z}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\beta}\right) \Longrightarrow \mathrm{Y}_{\mathrm{n}} \cdot \mathrm{Z}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha+\beta}\right)$ |
| 6 | $\mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right) \& \mathrm{Z}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\beta}\right) \Longrightarrow \Rightarrow \mathrm{Y}_{\mathrm{n}} \cdot \mathrm{Z}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\alpha+\beta}\right)$ |
| 7 | $\mathrm{Y}_{\mathrm{n}}=o_{\mathrm{o}}\left(\mathrm{n}^{\alpha}\right) \& \mathrm{Z}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\beta}\right) \& \beta \geq \alpha \Longrightarrow \mathrm{Y}_{\mathrm{n}}+\mathrm{Z}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\beta}\right)$ |
| 8 | $\mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right) \& \mathrm{Z}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\beta}\right) \& \beta \geq \alpha \Longrightarrow \mathrm{Y}_{\mathrm{n}}+\mathrm{Z}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\beta}\right)$ |
| 9 | $\mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right) \& \mathrm{Z}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\beta}\right) \& \beta>\alpha \Longrightarrow \mathrm{Y}_{\mathrm{n}}+\mathrm{Z}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\beta}\right)$ |
| 10 | $\mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right) \& \mathrm{Z}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\beta}\right) \& \beta<\alpha \Longrightarrow \mathrm{Y}_{\mathrm{n}}+\mathrm{Z}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right)$ |
| 11 | $\mathrm{Y}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right) \& \mathrm{Z}_{\mathrm{n}}=o_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right) \Longrightarrow \mathrm{Y}_{\mathrm{n}}+\mathrm{Z}_{\mathrm{n}}=O_{\mathrm{p}}\left(\mathrm{n}^{\alpha}\right)$ |

### 4.2. INDEPENDENT AND DEPENDENT RANDOM SEQUENCES

4.2.1. Consider a sequence of random variables $\mathrm{Y}_{1}, \mathrm{Y}_{2}, \mathrm{Y}_{3}, \ldots$. The joint distribution (CDF) of a finite subsequence $\left(Y_{1}, \ldots, Y_{n}\right)$, denoted $F_{1, \ldots, n}\left(y_{1}, \ldots, y_{n}\right)$, is defined as the probability of a state of Nature such that all of the inequalities $\mathrm{Y}_{1} \leq \mathrm{y}_{1}, \ldots, \mathrm{Y}_{\mathrm{n}} \leq \mathrm{y}_{\mathrm{n}}$ hold. The random variables in the sequence are mutually statistically independent if for every finite subsequence $\mathrm{Y}_{1} \ldots, \mathrm{Y}_{\mathrm{n}}$, the joint CDF factors:

$$
\mathrm{F}_{1, \ldots, \mathrm{n}}\left(\mathrm{y}_{1}, \ldots, \mathrm{y}_{\mathrm{n}}\right) \equiv \mathrm{F}_{1}\left(\mathrm{y}_{1}\right) \cdot \ldots \cdot \mathrm{F}_{\mathrm{n}}\left(\mathrm{y}_{\mathrm{n}}\right) .
$$

The variables are independent and identically distributed (i.i.d.) if in addition they have a common univariate $\operatorname{CDF~}_{1}(\mathrm{y})$. The case of i.i.d. random variables leads to the simplest theory of stochastic limits, and provides the foundation needed for much of basic econometrics. However, there are
many applications, particularly in analysis of economic time series, where i.i.d. assumptions are not plausible, and a limit theory is needed for dependent random variables. We will define two types of dependence, martingale and mixing, that will cover a variety of econometric time series applications and require a modest number of tools from probability theory. We have introduced a few of the needed tools in Chapter 3, notably the idea of information contained in $\sigma$-fields of events, with the evolution of information captured by refinements of these $\sigma$-fields, and the definitions of measurable functions, product $\sigma$-fields, and compatability conditions for probabilities defined on product spaces. There are treatments of more general forms of dependence than martingale or mixing, but these require a more comprehensive development of the theory of stochastic processes.
4.2.2. Consider a sequence of random variables $\mathrm{Y}_{\mathrm{k}}$ with k interpreted as an index of (discrete) time. One can think of k as the infinite sequence $\mathrm{k} \in \mathbf{K}=\{1,2, \ldots\}$, or as a doubly infinite sequence, extending back in time as well as forward, $\mathrm{k} \in \mathbf{K}=\{\ldots,-2,-1,0,1,2, \ldots\}$. The set of states of Nature can be defined as the product space $\mathbf{S}=\mathbf{X}_{i \in \mathbf{K}} \mathbb{R}$, or $\mathbf{S}=\mathbb{R}^{\mathbf{K}}$, where $\mathbb{R}$ is the real line, and the "complete information" $\sigma$-field of subsets of $\mathbf{S}$ defined as $\boldsymbol{F}_{\mathbf{K}}=\otimes_{\mathrm{i} \in \mathrm{K}} \boldsymbol{B}$, where $\boldsymbol{B}$ is the Borel $\sigma$-field of subsets of the real line; see 3.2. (The same apparatus, with $\mathbf{K}$ equal to the real line, can be used to consider continuous time. To avoid a variety of mathematical technicalities, we will not consider the continuous time case here.) Accumulation of information is described by a nondecreasing sequence of $\boldsymbol{\sigma}$-fields $\ldots \subseteq \boldsymbol{G}_{-1} \subseteq \boldsymbol{G}_{0} \subseteq \boldsymbol{G}_{1} \subseteq \boldsymbol{G}_{2} \subseteq \ldots$, with $\boldsymbol{G}_{\mathrm{t}}=\left(\bigotimes_{\mathrm{i} \leq \mathrm{t}} \boldsymbol{B}\right) \otimes\left(\bigotimes_{\mathrm{i}>\mathrm{t}}\{\boldsymbol{\varphi}, \mathbf{S}\}\right)$ capturing the idea that at time $t$ the future is unknown. The monotone sequence of $\sigma$-fields $\boldsymbol{G}_{v}, \mathrm{i}=\ldots,-1,0,1,2, \ldots$ is called a filtration. The sequence of random variables $\mathrm{Y}_{\mathrm{t}}$ is adapted to the filtration if $\mathrm{Y}_{\mathrm{t}}$ is measurable with respect to $\boldsymbol{G}_{\mathrm{t}}$ for each t . Some authors use the notation $\sigma\left(\ldots, \mathrm{Y}_{\mathrm{t}-2}, \mathrm{Y}_{\mathrm{t}-1}, \mathrm{Y}_{\mathrm{t}}\right)$ for $\boldsymbol{G}_{\mathrm{t}}$ to emphasize that it is the $\sigma$-field generated by the information contained in $Y_{s}$ for $s \leq t$. The sequence $\ldots, \mathrm{Y}_{-1}, \mathrm{Y}_{0}, \mathrm{Y}_{1}, \mathrm{Y}_{2}, \ldots$ adapted to $\boldsymbol{G}_{\mathrm{t}}$ for $\mathrm{k} \in \mathbf{K}$ is termed a stochastic process. One way of thinking of a stochastic process is to recall that random variables are functions of states of Nature, so that the process is a function $\mathrm{Y}: \mathrm{S} \times \mathrm{K} \rightarrow \mathbb{R}$. Then $\mathrm{Y}(\mathrm{s}, \mathrm{k})$ is the realization of the random variable in period k , $\mathrm{Y}(\mathrm{s}, \cdot)$ a realization or time-path of the stochastic process, and $\mathrm{Y}(\cdot, \mathrm{k})$ the random variable in period k . Note that there may be more than one sequence of $\sigma$-fields in operation for a particular process. These might correspond, for example, to the information available to different economic agents. We will need in particular the sequence of $\sigma$-fields $\boldsymbol{H}_{\mathrm{t}}=\sigma\left(\mathrm{Y}_{\mathrm{t}}, \mathrm{Y}_{\mathrm{t}+1}, \mathrm{Y}_{\mathrm{t}+2}, \ldots\right)$ adapted to the process from time $t$ forward; this is a nonincreasing sequence of $\sigma$-fields... $\supseteq \boldsymbol{H}_{\mathrm{t}-1} \supseteq \boldsymbol{H}_{\mathrm{t}} \supseteq \boldsymbol{H}_{\mathrm{t}+1} \supseteq \ldots$. Sometimes $\boldsymbol{G}_{\mathrm{t}}$ is termed the natural upward filtration, and $\boldsymbol{H}_{\mathrm{t}}$ the natural downward filtration.

Each subsequence $\left(\mathrm{Y}_{\mathrm{m}}, \ldots, \mathrm{Y}_{\mathrm{m}+\mathrm{n}}\right)$ of the stochastic process has a multivariate CDF $F_{m, \ldots m+n}\left(y_{m}, \ldots, y_{m+n}\right)$. It is said to be stationary if for each $n$, this CDF is the same for every m. A stationary process has the obvious property that moments such as means, variances, and covariances between random variables a fixed number of time periods apart are the same for all times m . Referring to 4.2.1, a sequence i.i.d. random variables is always stationary.
4.2.3. One circumstance that arises in some economic time series is that while the successive random variables are not independent, they have the property that their expectation, given history, is zero. Changes in stock market prices, for example, will have this property if the market is efficient, with arbitragers finding and bidding away any component of change that is predictable from history. A sequence of random variables $X_{t}$ adapted to $\boldsymbol{G}_{\mathrm{t}}$ is a martingale if almost surely $\mathbf{E}\left\{\mathrm{X}_{\mathrm{t}} \mid \boldsymbol{G}_{\mathrm{t}-1}\right\}=\mathrm{X}_{\mathrm{t}-1}$. If $\mathrm{X}_{\mathrm{t}}$ is a martingale, then $\mathrm{Y}_{\mathrm{t}}=\mathrm{X}_{\mathrm{t}}-\mathrm{X}_{\mathrm{t}-1}$ satisfies $\mathbf{E}\left\{\mathrm{Y}_{\mathrm{t}} \mid \boldsymbol{G}_{\mathrm{t}-1}\right)=0$, and is called a martingale difference (m.d.) sequence. Thus, stock price changes in an efficient market form a m.d. sequence. It is also useful to define a supermartingale (resp., submartingale) if almost surely $\mathbf{E}\left\{\mathrm{X}_{\mathrm{t}} \mid \boldsymbol{G}_{\mathrm{t}-1}\right) \leq \mathrm{X}_{\mathrm{t}-1}\left(\mathrm{resp},, \mathbf{E}\left\{\mathrm{X}_{\mathrm{t}} \mid \boldsymbol{G}_{\mathrm{t}-1}\right) \geq \mathrm{X}_{\mathrm{t}-1}\right.$ ). The following result, called the Kolmogorov maximal inequality, is a useful property of martingale difference sequences.

Theorem 4.5. If random variables $Y_{k}$ are have the property that $E\left(Y_{k} \mid Y_{1} .,,, . Y_{k-1}\right)=0$, or more technically the property that $\mathrm{Y}_{\mathrm{k}}$ adapted to $\sigma\left(\ldots, \mathrm{Y}_{\mathrm{k}-1}, \mathrm{Y}_{\mathrm{k}}\right)$ is a martingale difference sequence, and if $\mathbf{E} \mathrm{Y}_{\mathrm{k}}{ }^{2}=\sigma_{\mathrm{k}}{ }^{2}$, then $\mathrm{P}\left(\max _{1 \leq \mathrm{k} \leq \mathrm{n}}\left|\quad \sum_{i=1}^{k} \quad \mathrm{Y}_{\mathrm{i}}\right|>\varepsilon\right) \leq \sum_{i=1}^{n} \sigma_{\mathrm{i}}^{2} / \varepsilon^{2}$.

Proof: Let $\mathrm{S}_{\mathrm{k}}=\sum_{i=1}^{k} \quad \mathrm{Y}_{\mathrm{i}}$. Let $\mathrm{Z}_{\mathrm{k}}$ be a random variable that is one if $\mathrm{S}_{\mathrm{j}} \leq \boldsymbol{\varepsilon}$ for $\mathrm{j}<\mathrm{k}$ and $\mathrm{S}_{\mathrm{k}}>\varepsilon$, zero otherwise. Note that $\sum_{i=1}^{n} \mathrm{Z}_{\mathrm{i}} \leq 1$ and $\mathbf{E}\left(\sum_{i=1}^{n} \mathrm{Z}_{\mathrm{i}}\right)=\mathrm{P}\left(\max _{1 \leq \mathrm{k} \leq \mathrm{n}}\left|\sum_{i=1}^{k} \mathrm{Y}_{\mathrm{i}}\right|>\varepsilon\right)$. The variables $S_{k}$ and $Z_{k}$ depend only on $Y_{i}$ for $\mathrm{i} \leq \mathrm{k}$. Then $\mathrm{E}\left(\mathrm{S}_{\mathrm{n}}-\mathrm{S}_{\mathrm{k}} \mid \mathrm{S}_{\mathrm{k}}, \mathrm{Z}_{\mathrm{k}}\right)=0$. Hence

$$
\mathbf{E S}_{\mathrm{n}}{ }^{2} \geq \sum_{k=1}^{n} \quad \mathbf{E S}_{\mathrm{n}}{ }^{2} \cdot \mathrm{Z}_{\mathrm{k}}=\sum_{k=1}^{n} \quad \mathbf{E}\left[\mathrm{~S}_{\mathrm{k}}+\left(\mathrm{S}_{\mathrm{n}}-\mathrm{S}_{\mathrm{k}}\right)\right]^{2} \cdot \mathrm{Z}_{\mathrm{k}} \geq \sum_{k=1}^{n} \quad \mathbf{E S}_{\mathrm{k}}{ }^{2} \cdot \mathrm{Z}_{\mathrm{k}} \geq \boldsymbol{\varepsilon}^{2} \sum_{k=1}^{n} \quad \mathbf{E} Z_{\mathrm{k}} .
$$

4.2.4. As a practical matter, many economic time series exhibit correlation between different time periods, but these correlations dampen away as time differences increase. Bounds on correlations by themselves are typically not enough to give a satisfactory theory of stochastic limits, but a related idea is to postulate that the degree of statistical dependence between random variables approaches negligibility as the variables get further apart in time, because the influence of ancient history is buried in an avalance of new information (shocks). To formalize this, we introduce the concept of stochastic mixing. For a stochastic process $Y_{t}$, consider events $\mathbf{A} \in \boldsymbol{G}_{\mathrm{t}}$ and $\mathbf{B} \in \boldsymbol{H}_{t+s}$; then A draws only on information up through period tand $\mathbf{B}$ draws only on information from period $\mathrm{t}+\mathrm{s}$ on. The idea is that when $s$ is large, the information in $\mathbf{A}$ is too "stale" to be of much use in determining the probability of $\mathbf{B}$, and these events are nearly independent. Three definitions of mixing are given in the table below; they differ only in the manner in which they are normalized, but this changes their strength in terms of how broadly they hold and what their implications are. When the process is stationary, mixing depends only on time differences, not on time location.


There are links between the mixing conditions and bounds on correlations between events that are remote in time:
(1) Strict mixing $=\Rightarrow$ Uniform mixing $\Longrightarrow$ Strong mixing.
(2) (Serfling) If the $Y_{i}$ are uniform mixing with $\mathbf{E Y} Y_{i}=0$ and $\mathbf{E Y}{ }_{t}^{2}=\sigma_{t}^{2}<+\infty$, then $\left|\mathbf{E} Y_{\mathrm{t}} \mathrm{Y}_{\mathrm{t}+\mathrm{s}}\right| \leq 2 \varphi(\mathrm{~s})^{1 / 2} \sigma_{\mathrm{t}} \sigma_{\mathrm{t}+\mathrm{s}}$.
(3) (Ibragimov) If the $Y_{i}$ are strong mixing with $\mathbf{E} Y_{t}=0$ and $\mathbf{E}\left|Y_{t}\right|^{d}<+\infty$ for some $d>2$, then $\left|\mathbf{E Y}_{\mathrm{t}} \mathrm{Y}_{\mathrm{t}+\mathrm{s}}\right| \leq 8 \alpha(\mathrm{~s})^{1-2 / \mathrm{d}} \sigma_{\mathrm{t}} \sigma_{\mathrm{t}+\mathrm{s}}$.
(4) If there exists a sequence $\rho_{\mathrm{t}}$ with $\lim _{\mathrm{t}-\infty} \rho_{\mathrm{t}}=0$ such that $|\mathbf{E}(\mathbf{U}-\mathbf{E U})(\mathrm{W}-\mathbf{E W})| \leq$ $\rho_{\mathrm{t}}\left[\left(\mathbf{E}(\mathrm{U}-\mathbf{E U})^{2}\right)\left(\mathbf{E}(\mathrm{W}-\mathbf{E W})^{2}\right)\right]^{1 / 2}$ for all bounded continuous functions $\mathrm{U}=\mathrm{g}\left(\mathrm{Y}_{1}, \ldots, \mathrm{Y}_{\mathrm{t}}\right)$ and $\mathrm{W}=$ $h\left(Y_{t+n}, \ldots, Y_{t+n+m}\right)$ and all $t, n, m$, then the $Y_{t}$ are strict mixing.

An example gives an indication of the restrictions on a dependent stochastic process that produce strong mixing at a specified rate. First, suppose a stationary stochastic process $Y_{t}$ satisfies $Y_{t}=\rho Y_{t-1}$ $+Z_{t}$, with the $Z_{t}$ imdependent standard normal. Then, $\operatorname{var}\left(Y_{t}\right)=1 /\left(1-\rho^{2}\right)$ and $\operatorname{cov}\left(Y_{t+s}, Y_{t}\right)=\rho^{s} /\left(1-\rho^{2}\right)$, and one can show with a little analysis that $\left|P\left(Y_{t+s} \leq a, Y_{t} \leq b\right)-P\left(Y_{t+s} \leq a\right) \cdot P\left(Y_{t} \leq b\right)\right| \leq \rho^{5} / \pi\left(1-\rho^{2 s}\right)^{1 / 2}$. Hence, this process is strong mixing with a mixing coefficient that declines at a geometric rate. This is true more generally of processes that are formed by taking stationary linear transformations of independent processes. We return to this subject in the chapter on time series analysis.

### 4.3. LAWS OF LARGE NUMBERS

4.3.1. Consider a sequence of random variables $Y_{1}, Y_{2}, \ldots$ and a corresponding sequence of averages $\mathrm{X}_{\mathrm{n}}=n^{-1} \sum_{i=1}^{n} \quad \mathrm{Y}_{\mathrm{i}}$ for $\mathrm{n}=1,2, \ldots$. Laws of large numbers give conditions under which the averages $X_{n}$ converge to a constant, either in probability (weak laws, or WLLN) or almost surely (strong laws, or SLLN). Laws of large numbers give formal content to the intuition that sample averages are accurate analogs of population averages when the samples are large, and are essential to establishing that statistical estimators for many problems have the sensible property that with sufficient data they are likely to be close to the population values they are trying to estimate. In
econometrics, convergence in probability provided by a WLLN suffices for most purposes. However, the stronger result of almost sure convergence is occasionally useful, and is often attainable without additional assumptions.
4.3.2 Figure 4.4 lists a sequence of laws of large numbers. The case of independent identically distributed (i.i.d.) random variables yields the strongest result (Kolmogorov I). With additional conditions it is possible to get a laws of large numbers even for correlated variable provided the correlations of distant random variables approach zero sufficiently rapidly.

FIGURE 4.4. LAWS OF LARGE NUMBERS FOR $X_{\mathrm{n}}=n^{-1} \sum_{k=1}^{n} \mathrm{Y}_{\mathrm{k}}$

## WEAK LAWS (WLLN)

1 (Khinchine) If the $\mathrm{Y}_{\mathrm{k}}$ are i.i.d., and $\mathbf{E} \mathrm{Y}_{\mathrm{k}}=\mu$, then $X_{\mathrm{n}}{ }_{\mathrm{p}} \mu$
2 (Chebyshev) If the $\mathrm{Y}_{\mathrm{k}}$ are uncorrelated with $\mathbf{E} \mathrm{Y}_{\mathrm{k}}=\mu$ and $\mathbf{E}\left(\mathrm{Y}_{\mathrm{k}}-\mu\right)^{2}=\sigma_{\mathrm{k}}^{2}$ satisfying

$$
\sum_{k=1}^{\infty} \quad \sigma_{\mathrm{k}}^{2} / \mathrm{k}^{2}<+\infty, \text { then } \mathrm{X}_{\mathrm{n}} \overrightarrow{\mathrm{p}} \mu
$$

3 If the $\mathrm{Y}_{\mathrm{k}}$ have $\mathbf{E} \mathrm{Y}_{\mathrm{k}}=\mu, \mathbf{E}\left(\mathrm{Y}_{\mathrm{k}}-\mu\right)^{2} \equiv \sigma_{\mathrm{k}}^{2}$, and $\left|\mathbf{E}\left(\mathrm{Y}_{\mathrm{k}}-\mu\right)\left(\mathrm{Y}_{\mathrm{m}}-\mu\right)\right| \leq \rho_{\mathrm{km}} \sigma_{\mathrm{k}} \sigma_{\mathrm{m}}$ with

$$
\sum_{k=1}^{\infty} \quad \sigma_{\mathrm{k}}^{2} / \mathrm{k}^{3 / 2}<+\infty \text { and } \lim _{\mathrm{n} \rightarrow \infty} \frac{1}{n} \sum_{k=1}^{n} \sum_{m=1}^{n} \rho_{\mathrm{km}}<+\infty, \text { then } \mathrm{X}_{\mathrm{n}} \rightarrow_{\mathrm{p}} \mu
$$

## STRONG LAWS (SLLN)

1 (Kolmogorov I) If the $\mathrm{Y}_{\mathrm{k}}$ are i.i.d., and $\mathbf{E} \mathrm{Y}_{\mathrm{k}}=\mu$, then $\mathrm{X}_{\mathrm{n}} \rightarrow_{\text {as }} \mu$
2 (Kolmogorov II) If the $\mathrm{Y}_{\mathrm{k}}$ are independent, with $\mathbf{E} \mathrm{Y}_{\mathrm{k}}=\mu$, and $\mathbf{E}\left(\mathrm{Y}_{\mathrm{k}}-\mu\right)^{2}=\sigma_{\mathrm{k}}{ }^{2}$ satisfying

$$
\sum_{k=1}^{\infty} \quad \sigma_{\mathrm{k}}^{2} / \mathrm{k}^{2}<+\infty, \text { then } \mathrm{X}_{\mathrm{n}} \rightarrow_{\text {as }} \mu
$$

3 (Martingale) $\mathrm{Y}_{\mathrm{k}}$ adapted to $\sigma\left(\ldots, \mathrm{Y}_{\mathrm{k}-1}, \mathrm{Y}_{\mathrm{k}}\right)$ is a martingale difference sequence, $\mathbf{E} Y_{\mathrm{t}}^{2}=\sigma_{\mathrm{t}}^{2}$, and $\sum_{k=1}^{\infty} \quad \sigma_{\mathrm{k}}^{2} / \mathrm{k}^{2}<+\infty$, then $\mathrm{X}_{\mathrm{n}} \rightarrow_{\text {as }} 0$

4 (Serfling) If the $\mathrm{Y}_{\mathrm{k}}$ have $\mathbf{E} \mathrm{Y}_{\mathrm{k}}=\mu, \mathbf{E}\left(\mathrm{Y}_{\mathrm{k}}-\mu\right)^{2}=\sigma_{\mathrm{k}}{ }^{2}$, and $\left|\mathbf{E}\left(\mathrm{Y}_{\mathrm{k}}-\mu\right)\left(\mathrm{Y}_{\mathrm{m}}-\mu\right)\right| \leq \rho_{\mathrm{k}-\mathrm{m}} \sigma_{\mathrm{k}} \sigma_{\mathrm{m}}$,
with $\sum_{k=1}^{\infty}(\log k)^{2} \sigma_{k}^{2} / k^{2}<+\infty$ and $\sum_{k-1}^{\infty} \rho_{|k-m|}<+\infty$, then $X_{n} \rightarrow_{\text {as }} \mu$

To show why WLLN work, I outline proofs of the first three laws in Figure 4.4.
Theorem 4.6. (Khinchine) If the $Y_{k}$ are i.i.d., and $\mathbf{E} Y_{k}=\mu$, then $X_{n}{ }_{p} \mu$.
Proof: The argument shows that the characteristic function (c.f.) of $X_{n}$ converges pointwise to the c.f. for a constant random variable $\mu$. Let $\psi(t)$ be the c.f. of $Y_{1}$. Then $X_{n}$ has c.f. $\psi(t / n)^{n}$. Since $\mathbf{E Y} Y_{1}$ exists, $\psi$ has a Taylor's expansion $\psi(\mathrm{t})=1+\psi^{\prime}(\lambda \mathrm{t}) \mathrm{t}$, where $0<\lambda<1$ (see 3.5.12). Then $\psi(\mathrm{t} / \mathrm{n})^{\mathrm{n}}=$ $\left[1+(t / n) \psi^{\prime}(\lambda t / n)\right]^{n}$. But $\psi^{\prime}(\lambda t / n) \rightarrow \psi^{\prime}(0)=\mu$. A result from 2.1.10 states that if a sequence of scalars $\alpha_{n}$ has a limit, then $\left[1+\alpha_{n} / n\right]^{n} \rightarrow \exp \left(\lim \alpha_{n}\right)$. Then $\psi(t / n)^{n} \rightarrow e^{\mu t}$. But this is the c.f. of a constant random variable $\mu$, implying $X_{n}{ }_{d} \mu$, and hence $X_{n}{ }_{p} \mu$.

Theorem 4.7. (Chebyshev) If the $Y_{k}$ are uncorrelated with $\mathbf{E} Y_{k}=\mu$ and $\mathbf{E}\left(Y_{k}-\mu\right)^{2}=\sigma_{k}^{2}$ satisfying $\quad \sum_{k=1}^{\infty} \quad \sigma_{\mathrm{k}}^{2} / \mathrm{k}^{2}<+\infty$, then $\mathrm{X}_{\mathrm{n}} \overrightarrow{\mathrm{p}} \mu$.

Proof: One has $\mathbf{E}\left(\mathrm{X}_{\mathrm{n}}-\mu\right)^{2}=\sum_{k=1}^{n} \sigma_{n}^{2} / n^{2}$. Kronecker's Lemma (see 2.1.9) establishes that

$$
\sum_{k=1}^{\infty} \sigma_{k}^{2} / k^{2} \text { bounded implies } \mathbf{E}\left(X_{\mathrm{n}}-\mu\right)^{2} \rightarrow 0 . \text { Then Chebyshev's inequality implies } \mathrm{X}_{\mathrm{n}} \rightarrow_{\mathrm{p}} \mu
$$

The condition $\sum_{k=1}^{\infty} \sigma_{k}^{2} / k^{2}$ bounded in Theorem 4.7 is obviously satisfied if $\sigma_{\mathrm{k}}{ }^{2}$ is uniformly bounded, but is also satisfied if $\sigma_{\mathrm{k}}^{2}$ grows modestly with k ; e.g., it is sufficient to have $\sigma_{\mathrm{k}}{ }^{2}(\log \mathrm{~K}) / \mathrm{k}$ bounded.

Theorem 4.8. (WLLN 3) If the $\mathrm{Y}_{\mathrm{k}}$ have $\mathbf{E} \mathrm{Y}_{\mathrm{k}}=\mu, \mathbf{E}\left(\mathrm{Y}_{\mathrm{k}}-\mu\right)^{2} \equiv \sigma_{\mathrm{k}}^{2}$, and $\left|\mathbf{E}\left(\mathrm{Y}_{\mathrm{k}}-\mu\right)\left(\mathrm{Y}_{\mathrm{m}}-\mu\right)\right| \leq$ $\rho_{\mathrm{km}} \sigma_{\mathrm{k}} \sigma_{\mathrm{m}}$ with $\sum_{k=1}^{\infty} \sigma_{\mathrm{k}}^{2} / \mathrm{k}^{3 / 2}<+\infty$ and $\lim _{\mathrm{n} \rightarrow \infty} \frac{1}{n} \sum_{k=1}^{n} \sum_{m=1}^{n} \quad \rho_{\mathrm{km}}<+\infty$, then $\mathrm{X}_{\mathrm{n}} \overrightarrow{\mathrm{p}} \mu$

Proof: Using Chebyshev's inequality, it is sufficient to show that $\mathbf{E}\left(X_{n}-\mu\right)^{2}$ converges to zero. The Cauchy-Schwartz inequality (see 2.1.11) is applied first to establish

$$
\left(\frac{1}{n} \sum_{m=1}^{n} \sigma_{m} \rho_{k m}\right)^{2} \leq\left(\frac{1}{n} \sum_{m=1}^{n} \sigma_{m}^{2}\right)\left(\frac{1}{n} \sum_{m=1}^{n} \rho_{k m}^{2}\right)
$$

and then to establish that

$$
\mathbf{E}\left(\mathrm{X}_{\mathrm{n}}-\mu\right)^{2}=\frac{1}{n^{2}} \sum_{k=1}^{n} \sum_{m=1}^{n} \sigma_{k} \sigma_{m} \rho_{k m}=\frac{1}{n} \sum_{k=1}^{n} \sigma_{k}\left(\frac{1}{n} \sum_{m=1}^{n} \sigma_{m} \rho_{k m}\right)
$$

$$
\begin{aligned}
\leq & \left(\frac{1}{n} \sum_{k=1}^{n} \sigma_{k}^{2}\right)^{1 / 2}\left[\frac{1}{n} \sum_{k=1}^{n}\left(\frac{1}{n} \sum_{m=1}^{n} \sigma_{m} \rho_{k m}\right)^{2}\right]^{1 / 2} \leq\left(\frac{1}{n} \sum_{k=1}^{n} \sigma_{k}^{2}\right)^{1 / 2}\left[\left(\frac{1}{n} \sum_{m=1}^{n} \sigma_{m}^{2}\right)\left(\frac{1}{n^{2}} \sum_{k=1}^{n} \sum_{m=1}^{n} \rho_{k m}^{2}\right)\right]^{1 / 2} \\
& =\left(\frac{1}{n} \sum_{k=1}^{n} \sigma_{k}^{2}\right)\left(\frac{1}{n^{2}} \sum_{k=1}^{n} \sum_{m=1}^{n} \rho_{k m}^{2}\right)^{1 / 2}=\left(\frac{1}{n^{3 / 2}} \sum_{k=1}^{n} \sigma_{k}^{2}\right)\left(\frac{1}{n} \sum_{k=1}^{n} \sum_{m=1}^{n} \rho_{k m}^{2}\right)^{1 / 2} .
\end{aligned}
$$

The last form and Kronecker's lemma (2.1.11) give the result.
The conditions for this result are obviously met if the $\sigma_{\mathrm{k}}{ }^{2}$ are uniformly bounded and the correlation coefficients decline at a sufficient rate with the distance between observations; examples are geometric decline with $\rho_{\mathrm{km}}$ bounded by a multiple of $\lambda^{|k-m|}$ for some $\lambda<1$ and an arithmetic decline with $\rho_{\mathrm{km}}$ bounded by a multiple of $|\mathrm{k}-\mathrm{m}|^{-1}$.

The Kolmogorov SLLN 1 is a better result than the Kinchine WLLN, yielding a stronger conclusion from the same assumptions. Similarly, the Kolmogorov SLLN 2 is a better result than the Chebyshev WLLN. Proofs of these theorems can be found in C. R. Rao (1973), p. 114-115. The Serfling SLLN 4 is broadly comparable to WLLN 3, but Serfling gets the stronger almost sure conclusion with somewhat stronger assumptions on the correlations and somewhat weaker assumptions on the variances. If variances are uniformly bounded and correlation coefficients decline at least at a rate inversely proportional to the square of the time difference, this sufficient for either the WLLN 3 or SLLN 4 assumptions.

The SLLN 3 in the table applies to martingale difference sequences, and shows that Kolmogorov II actually holds for m.d. sequences.

Theorem 4.9. If $Y_{t}$ adapted to $\sigma\left(\ldots, Y_{k-1}, Y_{k}\right)$ is a martingale difference sequence with $\mathbf{E} Y_{t}^{2}=\sigma_{t}^{2}$ and $\sum_{k=1}^{\infty} \quad \sigma_{\mathrm{k}}^{2} / \mathrm{k}^{2}<+\infty$, then $\mathrm{X}_{\mathrm{n}} \rightarrow_{\text {as }} 0$.

Proof: The theorem is stated and proved by J. Davidson (1994), p. 314. To give an idea why SLLN work, I will give a simplified proof when the assumption $\sum_{k=1}^{\infty} \sigma_{\mathrm{k}}^{2} / \mathrm{k}^{2}<+\infty$ is strengthened to $\sum_{k=1}^{\infty} \sigma_{\mathrm{k}}^{2} / \mathrm{k}^{3 / 2}<+\infty$. Either assumption handles the case of constant variances with room to spare. Kolmogorov's maximal inequality (Theorem 4.5) with $\mathrm{n}=(\mathrm{m}+1)^{2}$ and $\varepsilon=\delta \mathrm{m}^{2}$ implies that

$$
\mathrm{P}\left(\max _{\mathrm{m}^{2} \leq \mathrm{k} s(\mathrm{~m}+1)^{2}}\left|\mathrm{X}_{\mathrm{k}}\right|>\delta\right) \leq \mathrm{P}\left(\max _{1 \leq \mathrm{k} \leq \mathrm{n}}\left|\quad \sum_{i=1}^{k} \quad \mathrm{Y}_{\mathrm{i}}\right|>\delta \mathrm{m}^{2}\right) \leq \sum_{i=1}^{(m+1)^{2}} \sigma_{\mathrm{i}}^{2} / \delta^{2} \mathrm{~m}^{4}
$$

The sum over $m$ of the right-hand-side of this inequality satisfies

$$
\sum_{m=1}^{\infty} \sum_{i=1}^{(m+1)^{2}} \quad \sigma_{\mathrm{i}}^{2} / \delta^{2} \mathrm{~m}^{4}=\sum_{i=1}^{\infty} \sum_{m \geq i}^{\infty}{ }^{1 / 2} . \sigma_{\mathrm{i}}^{2} / \delta^{2} \mathrm{~m}^{4} \leq 36 \sum_{i=1}^{\infty} \quad \sigma_{\mathrm{i}}^{2} / \mathrm{i}^{3 / 2} \delta^{2}
$$

Then $\sum_{m=1}^{\infty} P\left(\sup _{\mathrm{k}}\left|\mathrm{X}_{\mathrm{k}}\right|>\delta\right) \leq 36 \sum_{i=1}^{\infty} \quad \sigma_{\mathrm{i}}^{2} / \mathrm{i}^{3 / 2} \delta^{2}<+\infty$. Theorem 4.2 gives the result.

### 4.4. CENTRAL LIMIT THEOREMS

4.4.1. Consider a sequence of random variables $Y_{1}, \ldots, Y_{n}$ with zero means, and the associated sequence of scaled averages $\mathrm{Z}_{\mathrm{n}}=\mathrm{n}^{-1 / 2} \quad \sum_{i=1}^{n} \quad \mathrm{Y}_{\mathrm{i}}$. Central limit theorems (CLT) are concerned with conditions under which the $Z_{n}$, or variants with more generalized scaling, converge in distribution to a normal random variable $Z_{0}$. I will present several basic CLT, prove the simplest, and discuss the remainder. These results are summarized in Figure 4.5.

The most straighforward CLT is obtained for independent and identically distributed (i.i.d.) random variables, and requires only that the random variables have a finite variance. Note that the finite variance assumption is an additional condition needed for the CLT that was not needed for the SLLN for i.i.d. variables.

Theorem 4.10. (Lindeberg-Levy) If random variables $Y_{k}$ are i.i.d. with mean zero and finite positive variance $\sigma^{2}$, then $\mathrm{Z}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{Z}_{\mathrm{o}} \sim N\left(0, \sigma^{2}\right)$.

Proof: The approach is to show that the characteristic function of $Z_{n}$ converges for each argument to the characteristic function of a normal. The CLT then follows from the limit properties of characteristic functions (see 3.5.12). Let $\psi(\mathrm{t})$ be the cf of $\mathrm{Y}_{1}$. Then $\mathrm{Z}_{\mathrm{n}}$ has $\mathrm{cf} \psi\left(\mathrm{t} \cdot \mathrm{n}^{-1 / 2}\right)^{\mathrm{n}}$. Since $\mathbf{E Y} Y_{1}$ $=0$ and $\mathbf{E Y} 1_{1}{ }^{2}=\sigma^{2}, \psi(\mathrm{t})$ has a Taylor's expansion $\psi(\mathrm{t})=\left[1+\psi^{\prime \prime}(\lambda \mathrm{t}) \mathrm{t}^{2} / 2\right]$, where $0<\lambda<1$ and $\psi^{\prime \prime}$ is continuous with $\psi^{\prime \prime}(0)=-\sigma^{2}$. Then $\psi\left(t \cdot n^{-1 / 2}\right)^{n}=\left[1+\psi^{\prime \prime}\left(\lambda t \cdot n^{-1 / 2}\right) t^{2} / 2 n\right]^{n}$. Then the limit result 2.1.10 gives $\lim _{n \rightarrow \infty}\left[1+\psi^{\prime \prime}\left(\lambda t \cdot n^{-1 / 2}\right) t^{2} / 2 n\right]^{n}=\exp \left(-\sigma^{2} t^{2} / 2\right)$. Thus, the $c f$ of $Z_{n}$ converges for each $t$ to the $c f$ of $Z_{o} \sim N\left(0, \sigma^{2}\right)$.
4.4.2. When the variables are independent but not identically distributed, an additional bound on the behavior of tails of the distributions of the random variables, called the Lindeberg condition, is needed. This condition ensures that sources of relatively large deviations are spread fairly evenly through the series, and not concentrated in a limited number of observations. The Lindeberg condition can be difficult to interpret and check, but there are a number of sufficient conditions that are useful in applications. The main result, stated next, allows more general scaling than by $\mathrm{n}^{-1 / 2}$.

FIGURE 4.5. CENTRAL LIMIT THEOREMS FOR $Z_{\mathbf{n}}=\mathbf{n}^{-1 / 2} \sum_{i=1}^{n} \mathbf{Y}_{\mathbf{i}}$
1 (Lindeberg-Levy) $\mathrm{Y}_{\mathrm{k}}$ i.i.d., $\mathbf{E Y} \mathrm{K}_{\mathrm{k}}=0, \mathbf{E} \mathrm{Y}_{\mathrm{k}}{ }^{2}=\sigma^{2}$ positive and finite $\Longrightarrow \Rightarrow \mathrm{Z}_{\mathrm{n}}{ }_{\mathrm{d}} \mathrm{Z}_{\mathrm{o}} \sim N\left(0, \sigma^{2}\right)$

2 (Lindeberg-Feller) If $\mathrm{Y}_{\mathrm{k}}$ independent, $\mathbf{E} \mathrm{Y}_{\mathrm{k}}=0, \mathbf{E Y}_{\mathrm{k}}{ }^{2}=\sigma_{\mathrm{k}}{ }^{2} \in(0,+\infty), \mathrm{c}_{\mathrm{n}}{ }^{2}=\sum_{k=1}^{n} \quad \sigma_{\mathrm{k}}{ }^{2}$, then $\mathrm{c}_{\mathrm{n}}{ }^{2} \rightarrow+\infty, \lim _{\mathrm{n} \rightarrow \infty} \max _{1 \leq \mathrm{k} \leqslant \mathrm{n}} \sigma_{\mathrm{k}} / \mathrm{c}_{\mathrm{n}}=0$, and $\mathrm{U}_{\mathrm{n}}=\sum_{k=1}^{n} \quad \mathrm{Y}_{\mathrm{k}} / \mathrm{c}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{U}_{\mathrm{o}} \sim N(0,1)$ if and only if the

Lindeberg condition holds: for each $\varepsilon>0, \sum_{k=1}^{n} \quad \mathbf{E} \mathrm{Y}_{\mathrm{k}}{ }^{2} \cdot \mathbf{1}\left(\left|\mathrm{Y}_{\mathrm{k}}\right|>\varepsilon \mathrm{c}_{\mathrm{n}}\right) / \mathrm{c}_{\mathrm{n}}{ }^{2} \rightarrow 0$

3
If $\mathrm{Y}_{\mathrm{k}}$ independent, $\mathbf{E} \mathrm{Y}_{\mathrm{k}}=0, \mathbf{E Y}_{\mathrm{k}}{ }^{2}=\sigma_{\mathrm{k}}{ }^{2} \in(0,+\infty), \mathrm{c}_{\mathrm{n}}{ }^{2}=\sum_{k=1}^{n} \quad \sigma_{\mathrm{k}}{ }^{2}$ have $\mathrm{c}_{\mathrm{n}}{ }^{2} \rightarrow+\infty$ and $\lim _{n \rightarrow \infty} \max _{1 \leq k \leq n} \sigma_{k} / c_{n}=0$, then each of the following conditions is sufficient for the Lindeberg condition:
(i) For some $\mathrm{r}>2, \quad \sum_{k=1}^{n} \quad \mathbf{E}\left|\mathrm{Y}_{\mathrm{k}}\right|^{\mathrm{r}} / \mathrm{c}_{\mathrm{n}}{ }^{\mathrm{r}} \rightarrow 0$.
(ii) (Liapunov) For some $\mathrm{r}>2, \mathbf{E}\left|\mathrm{Y}_{\mathrm{k}} / \sigma_{\mathrm{k}}\right|^{\mathrm{r}}$ is bounded uniformly for all n .
(iii) For some $\mathrm{r}>2, \mathbf{E}\left|\mathrm{Y}_{\mathrm{k}}\right|^{\mathrm{r}}$ is bounded, and $\mathrm{c}_{\mathrm{k}}^{2} / \mathrm{k}$ is bounded positive, uniformly for all k .
$4 \quad \mathrm{Y}_{\mathrm{k}}$ a martingale difference sequence adapted to $\sigma\left(\ldots, \mathrm{Y}_{\mathrm{k}-1}, \mathrm{Y}_{\mathrm{k}}\right)$ with $\left|\mathrm{Y}_{\mathrm{k}}\right|<\mathrm{M}$ for all t and $\mathbf{E} Y_{\mathrm{k}}{ }^{2}=\sigma_{\mathrm{k}}{ }^{2}$ satisfying $n^{-1} \sum_{k=1}^{n} \quad \sigma_{\mathrm{k}}{ }^{2} \rightarrow \sigma_{\mathrm{o}}{ }^{2}>0 \Longrightarrow \mathrm{Z}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{Z}_{\mathrm{o}} \sim N\left(0, \sigma_{\mathrm{o}}{ }^{2}\right)$

5 (Ibragimov-Linnik) $\mathrm{Y}_{\mathrm{k}}$ stationary and strong mixing with $\mathbf{E} \mathrm{Y}_{\mathrm{k}}=0, \mathbf{E} \mathrm{Y}_{\mathrm{k}}{ }^{2}=\sigma^{2} \in(0,+\infty)$, $\mathbf{E} Y_{k+s} Y_{k}=\sigma^{2} \rho_{s}$, and for some $\mathrm{r}>2, \mathbf{E}\left|\mathrm{Y}_{\mathrm{n}}\right|^{\mathrm{r}}<+\infty$ and $\sum_{k=1}^{\infty} \alpha(\mathrm{k})^{1-2 / \mathrm{r}}<+\infty \Longrightarrow \Rightarrow \sum_{s=1}^{\infty}\left|\rho_{s}\right|$ $<+\infty$ and $\mathrm{Z}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{Z}_{\mathrm{o}} \sim N\left(0, \sigma^{2}\left(1+2 \sum_{s=1}^{\infty} \quad \rho_{\mathrm{s}}\right)\right)$

Theorem 4.11. (Lindeberg-Feller) Suppose random variables $Y_{k}$ are independent with mean zero and positive finite variances $\sigma_{\mathrm{k}}{ }^{2}$. Define $\mathrm{c}_{\mathrm{n}}{ }^{2}=\sum_{k=1}^{n} \quad \sigma_{\mathrm{k}}{ }^{2}$ and $\mathrm{U}_{\mathrm{n}}=\sum_{k=1}^{n} \quad \mathrm{Y}_{\mathrm{k}} / \mathrm{c}_{\mathrm{n}}$. Then $\mathrm{c}_{\mathrm{n}}{ }^{2} \rightarrow \infty$, $\lim _{\mathrm{n} \rightarrow \infty} \max _{1 \leq \mathrm{k} \leq \mathrm{n}} \sigma_{\mathrm{k}} / \mathrm{c}_{\mathrm{n}}=0$, and $\mathrm{U}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{U}_{\mathrm{o}} \sim N(0,1)$ if and only if the $\mathrm{Y}_{\mathrm{k}}$ satisfy the Lindeberg condition that for $\varepsilon>0, \lim _{\mathrm{n} \rightarrow \infty} \quad \sum_{i=1}^{n} \quad \mathbf{E} \mathrm{Y}_{\mathrm{k}}{ }^{2} \cdot \mathbf{1}\left(\left|\mathrm{Y}_{\mathrm{k}}\right|>\varepsilon \mathrm{c}_{\mathrm{n}}\right) / \mathrm{c}_{\mathrm{n}}{ }^{2}=0$.

A proof of Theorem 4.11 can be found, for example, in P. Billingsley (1986), p. 369-375. It involves an analysis of the characteristic functions, with detailed analysis of the remainder terms in their Taylor's expansion. To understand the theorem, it is useful to first specialize it to the case that the $\sigma_{\mathrm{k}}^{2}$ are all the same. Then $\mathrm{c}_{\mathrm{n}}{ }^{2}=\mathrm{n} \sigma_{1}^{2}$, the conditions $\mathrm{c}_{\mathrm{n}}{ }^{2} \rightarrow \infty$ and $\lim _{\mathrm{n} \rightarrow \infty} \max _{1 \leq \mathrm{k} \leq \mathrm{n}} \sigma_{\mathrm{k}} / \mathrm{c}_{\mathrm{n}}=0$ are met automatically, and in the terminology at the start of this section, $\mathrm{U}_{\mathrm{n}}=\mathrm{Z}_{\mathrm{n}} / \sigma_{1}$. The theorem then says $\mathrm{U}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{U}_{\mathrm{o}} \sim N(0,1)$ if and only if the sample average of $\mathbf{E} \mathrm{Y}_{\mathrm{k}}{ }^{2} \cdot \mathbf{1}\left(\left|\mathrm{Y}_{\mathrm{k}}\right|>\varepsilon \mathrm{n}^{1 / 2}\right)$ converges to zero for each $\varepsilon>0$. The last condition limits the possibility that the deviations in a single random variable could be as large in magnitude as the sum, so that the shape of the distribution of this variable makes a significant contribution to the shape of the distribution of the sum. An example shows how the Lindeberg condition bites. Consider independent random variables $\mathrm{Y}_{\mathrm{k}}$ that equal $\pm \mathrm{k}^{\mathrm{r}}$ with probability $1 / 2 \mathrm{k}^{2 \mathrm{r}}$, and zero otherwise, where r is a positive scalar. The $\mathrm{Y}_{\mathrm{k}}$ have mean zero and variance one, and $\mathbf{1}\left(\left|\mathrm{Y}_{\mathrm{k}}\right|>\varepsilon \boldsymbol{n}^{1 / 2}\right)=1$ if $\mathrm{k}^{\mathrm{r}}>\boldsymbol{\varepsilon} \mathbf{n}^{1 / 2}$, implying $n^{-1} \sum_{i=1}^{n} \quad \mathbf{E} \mathrm{Y}_{\mathrm{k}}{ }^{2} \cdot \mathbf{1}\left(\left|\mathrm{Y}_{\mathrm{k}}\right|>\varepsilon \mathbf{n}^{1 / 2}\right)=\max \left(0,1-\varepsilon^{1 / \mathrm{r}} \mathbf{n}^{(1-2 \mathrm{r} / 2 \mathrm{r}}\right)$.

This converges to zero, so the Lindeberg condition is satisfied iff $r<1 / 2$. Thus, the tails of the sequence of random variables cannot "fatten" too rapidly.

The Lindeberg condition allows the variances of the $\mathrm{Y}_{\mathrm{k}}$ to vary within limits. For example, the variables $Y_{k}= \pm 2^{k}$ with probability $1 / 2$ have $\sigma_{n} / c_{n}$ bounded positive, so that the variances grow too rapidly and the condition fails. The variables $Y_{k}= \pm 2^{-k}$ with probability $1 / 2$ have $c_{n}$ bounded, so that $\sigma_{1} / c_{n}$ is bounded positive, the variances shrink too rapidly, and the condition fails. The next result gives some easily checked sufficient conditions for the Lindeberg condition.

Theorem 4.12. Suppose random variables $Y_{k}$ are independent with mean zero and positive finite variances $\sigma_{\mathrm{k}}^{2}$ that satisfy $\mathrm{c}_{\mathrm{n}}^{2}=\sum_{k=1}^{n} \quad \sigma_{\mathrm{k}}{ }^{2} \rightarrow \infty$ and $\lim _{\mathrm{n} \rightarrow \infty} \max _{1 \leq \mathrm{k} \leq \mathrm{n}} \sigma_{\mathrm{k}} / \mathrm{c}_{\mathrm{n}}=0$. Then, each of the following conditions is sufficient for the Lindeberg condition to hold:
(i) For some $\mathrm{r}>2, \quad \sum_{k=1}^{n} \quad \mathbf{E}\left|\mathrm{Y}_{\mathrm{k}}\right|^{\mathrm{r}} / \mathrm{c}_{\mathrm{n}}{ }^{\mathrm{r}} \rightarrow 0$.
(ii) (Liapunov) For some $\mathrm{r}>2, \mathbf{E} \mid \mathrm{Y}_{\mathrm{k}} / \sigma_{\mathrm{k}}{ }^{\mathrm{r}}$ is bounded uniformly for all n .
(iii) For some $\mathrm{r}>2, \mathbf{E} \mid \mathrm{Y}_{\mathrm{k}}{ }^{\mathrm{l}}$ is bounded, and $\mathrm{c}_{\mathrm{k}}^{2} / \mathrm{k}$ is bounded positive, uniformly for all k .

Proof: To show that (i) implies the Lindeberg condition, write

$$
\sum_{k=1}^{n} \quad \mathbf{E} \mathrm{Y}_{\mathrm{k}}^{2} \cdot \mathbf{1}\left(\left|\mathrm{Y}_{\mathrm{k}}\right|>\varepsilon \mathrm{c}_{\mathrm{n}}\right) / \mathrm{c}_{\mathrm{n}}{ }^{2} \leq\left(\varepsilon \mathrm{c}_{\mathrm{n}}\right)^{2-\mathrm{r}} \quad \sum_{k=1}^{n} \quad \mathbf{E}\left|\mathrm{Y}_{\mathrm{k}}\right|^{\mathrm{r}} \cdot \mathbf{1}\left(\left|\mathrm{Y}_{\mathrm{k}}\right|>\varepsilon \mathrm{c}_{\mathrm{n}}\right) / \mathrm{c}_{\mathrm{n}}{ }^{2} \leq \boldsymbol{\varepsilon}^{2-\mathrm{r}} \quad \sum_{k=1}^{n} \quad \mathbf{E}\left|\mathrm{Y}_{\mathrm{k}} / \mathrm{c}_{\mathrm{n}}\right|^{\mathrm{r}} .
$$

For (ii), the middle expression in the string of inequalities above satisfies

$$
\begin{gathered}
\left(\varepsilon \mathrm{c}_{\mathrm{n}}\right)^{2-\mathrm{r}} \quad \sum_{k=1}^{n} \quad \mathbf{E} \mid \mathrm{Y}_{\mathrm{k}}{ }^{\mathrm{r}} \cdot \mathbf{1}\left(\left|\mathrm{Y}_{\mathrm{k}}\right|>\mathrm{c}_{\mathrm{n}}\right) / \mathrm{c}_{\mathrm{n}}^{2} \leq \boldsymbol{\varepsilon}^{2-\mathrm{r}}\left(\max _{\mathrm{k} \leq \mathrm{n}} \mathbf{E}\left|\mathrm{Y}_{\mathrm{k}} / \sigma_{\mathrm{k}}\right|^{\mathrm{r}} \cdot \sum_{k=1}^{n} \quad \sigma_{\mathrm{k}}^{\mathrm{r}} / \mathrm{c}_{\mathrm{n}}^{\mathrm{r}}\right. \\
\\
\leq \boldsymbol{\varepsilon}^{2-\mathrm{r}}\left(\max _{\mathrm{k} \leq \mathrm{n}} \mathbf{E} \mid \mathrm{Y}_{\mathrm{k}} / \sigma_{\mathrm{k}}{ }^{\mathrm{r}}\right) \cdot \sum_{k=1}^{n} \quad\left(\sigma_{\mathrm{k}}^{2} / \mathrm{c}_{\mathrm{n}}^{2}\right) \cdot\left(\max _{\mathrm{k} \leq \mathrm{n}}\left(\sigma_{\mathrm{k}} / \mathrm{c}_{\mathrm{n}}\right)^{\mathrm{r}-2}\right),
\end{gathered}
$$

and the assumptions ensure that $\max _{\mathrm{k} \leq \mathrm{n}} \mathbf{E} \mid \mathrm{Y}_{\mathrm{k}} / \sigma_{\mathrm{k}}{ }^{\mathrm{r}}$ is bounded and $\max _{\mathrm{k} \leq \mathrm{n}}\left(\sigma_{\mathrm{k}} / \mathrm{c}_{\mathrm{n}}\right)^{\mathrm{r}-2} \rightarrow 0$.
Finally, if (iii), then continuing the first string of inequalities,

$$
\sum_{i=1}^{n} \quad \mathbf{E}\left|\mathrm{Y}_{\mathrm{k}}\right|^{\mathrm{r}} / \mathrm{c}_{\mathrm{n}}^{\mathrm{r}} \leq \mathrm{c}_{\mathrm{n}}{ }^{2-\mathrm{r}} \mathrm{n} \cdot\left(\sup _{\mathrm{k}} \mathbf{E} \mid \mathrm{Y}_{\mathrm{k}}{ }^{\mathrm{l}}\right) / \mathrm{n} \cdot\left(\mathrm{inf}_{\mathrm{n}} \mathrm{c}_{\mathrm{n}}{ }^{2} / \mathrm{n}\right),
$$

and the right-hand-side is proportional to $\mathrm{c}_{\mathrm{n}}{ }^{2-\mathrm{r}}$, which goes to zero.
4.4.3. The following theorem establishes a CLT for the scaled sum $Z_{n}=n^{-1 / 2} \quad \sum_{i=1}^{n} \quad Y_{i}$ of martingale differences; or $\mathrm{Z}_{\mathrm{n}}=\mathrm{n}^{-1 / 2}\left(\mathrm{X}_{\mathrm{n}}-\mathrm{X}_{\mathrm{o}}\right)$. The uniform boundedness assumption in this theorem is a strong restriction, but it can be relaxed to a Lindeberg condition or to a "uniform integratability" condition; see P. Billingsley (1984), p. 498-501, or J. Davidson (1994), p. 385. Martingale differences can display dependence that corresponds to important economic applications, such as conditional variances that depend systematically on history.

Theorem 4.13. Suppose $Y_{k}$ is a martingale difference adapted to $\sigma\left(\ldots, Y_{k-1}, Y_{k}\right)$, and $Y_{k}$ satisfies a uniform bound $\left|\mathrm{Y}_{\mathrm{k}}\right|<\mathrm{M}$. Let $\mathbf{E} \mathrm{Y}_{\mathrm{k}}{ }^{2}=\sigma_{\mathrm{k}}{ }^{2}$, and assume that $n^{-1} \sum_{k=1}^{n} \quad \sigma_{\mathrm{k}}{ }^{2} \rightarrow \sigma_{\mathrm{o}}^{2}>0$. Then $\mathrm{Z}_{\mathrm{n}} \overrightarrow{\mathrm{d}}$ $\mathrm{Z}_{\mathrm{o}} \sim N\left(0, \sigma_{\mathrm{o}}{ }^{2}\right)$.
4.4.4. Intuitively, the CLT results that hold for independent or martingale difference random variables should continue to hold if the degree of dependence between variables is negligible. The following theorem from I. Ibragimov and Y. Linnik, 1971, gives a CLT for stationary strong mixing processes. This result will cover a variety of economic applications, including stationary linear transformations of independent processes like the one given in the last example.

Theorem 4.14. (Ibragimov-Linnik) Suppose $Y_{k}$ is stationary and strong mixing with mean zero, variance $\sigma^{2}$, and covariances $\mathbf{E} \mathrm{Y}_{\mathrm{k}+\mathrm{s}} \mathrm{Y}_{\mathrm{k}}=\sigma^{2} \rho_{\mathrm{s}}$. Suppose that for some $\mathrm{r}>2, \mathbf{E}\left|\mathrm{Y}_{\mathrm{n}}\right|^{\mathrm{r}}<+\infty$ and

$$
\sum_{k=1}^{\infty} \alpha(\mathrm{k})^{1-2 / \mathrm{s}}<+\infty . \text { Then, } \quad \sum_{s=1}^{\infty} \quad\left|\rho_{\mathrm{s}}\right|<+\infty \text { and } \mathrm{Z}_{\mathrm{n}} \rightarrow_{\mathrm{d}} \mathrm{Z}_{\mathrm{o}} \sim N\left(0, \sigma^{2}\left(1+2 \sum_{s=1}^{\infty} \quad \rho_{\mathrm{s}}\right)\right) .
$$

### 4.5. EXTENSIONS OF LIMIT THEOREMS

4.5.1. Limit theorems can be extended in several directions: (1) obtaining results for "triangular arrays" that include weighted sums of random variables, (2) sharpening the rate of convergence to the limit for "well-behaved" random variables, and (3) establishing "uniform" laws that apply to random functions. In addition, there are a variety of alternatives to the cases given above where independence assumptions are relaxed. The first extension gives limit theorems for random variables weighted by other (non-random) variables, a situation that occurs often in econometrics. The second extension provides tools that allow us to bound the probability of large deviations of random sums. This is of direct interest as a sharper version of a Chebychev-type inequality, and also useful in obtaining further results. To introduce uniform laws, first define a random function (or stochastic process) $\mathrm{y}=\mathrm{Y}(\theta, \mathrm{s})$ that maps a state of Nature s and a real variable (or vector of variables) $\theta$ into the real line. This may also be written, suppressing the dependence on s, as $\mathrm{Y}(\theta)$. Note that $\mathrm{Y}(\cdot, \mathrm{w})$ is a realization of the random function, and is itself an ordinary non-random function of $\theta$. For each value of $\theta, \mathrm{Y}(\theta, \cdot)$ is an ordinary random variable. A uniform law is one that bounds sums of random functions uniformly for all arguments $\theta$. For example, a uniform WLLN would say $\lim _{n \rightarrow \infty}$ $\mathrm{P}\left(\sup _{\theta}\left|n^{-1} \sum_{i=1}^{n} \mathrm{Y}_{\mathrm{i}}(\theta, \cdot)\right|>\varepsilon\right)=0$. Uniform laws play an important role in establishing the properties of statistical estimators that are nonlinear functions of the data, such as maximum likelihood estimates.
4.5.2 Consider a doubly indexed array of constants $\mathrm{a}_{\mathrm{in}}$ defined for $1 \leq \mathrm{i} \leq \mathrm{n}$ and $\mathrm{n}=1,2, \ldots$, and weighted sums of the form $X_{n}=\sum_{i=1}^{n} a_{i n} Y_{i}$. If the $Y_{i}$ are i.i.d., what are the limiting properties of $X_{n}$ ? We next give a WLLN and a CLT for weighted sums. The way arrays like $a_{i n}$ typically arise is that there are some weighting constants $\mathrm{c}_{\mathrm{i}}$, and either $\mathrm{a}_{\mathrm{in}}=\mathrm{c}_{\mathrm{i}} / \sum_{i=1}^{n} \quad \mathrm{c}_{\mathrm{j}}$ or $\mathrm{a}_{\mathrm{in}}=\mathrm{c}_{\mathrm{i}} /\left[\sum_{i=1}^{n} \mathrm{c}_{\mathrm{j}}\right]^{1 / 2}$. If $\mathrm{c}_{\mathrm{i}}=1$ for all i , then $\mathrm{a}_{\mathrm{in}}=\mathrm{n}^{-1}$ or $\mathrm{n}^{-1 / 2}$, respectively, leading to the standard scaling in limit theorems.

Theorem 4.15. Assume random variables $\mathrm{Y}_{\mathrm{i}}$ are independently identically distributed with mean zero. If an array $\mathrm{a}_{\mathrm{in}}$ satisfies $\lim _{\mathrm{n} \rightarrow \infty} \sum_{i=1}^{n}\left|\mathrm{a}_{\mathrm{jn}}\right|=0$ and $\lim _{\mathrm{n} \rightarrow \infty} \max _{\mathrm{j} \leq \mathrm{n}}\left|\mathrm{a}_{\mathrm{jn}}\right|=0$, then $\mathrm{X}_{\mathrm{n}} \rightarrow_{\mathrm{p}} 0$.

Proof: This is a weighted version of Khinchine's WLLN, and is proved in the same way. Let $\zeta(\mathrm{t})$ be the second characteristic function of $Y_{1}$. From the properties of characteristic functions we have $\zeta^{\prime}(0)=0$ and a Taylor's expansion $\zeta(\mathrm{t})=\mathrm{t} \cdot \zeta^{\prime}(\lambda \mathrm{t})$ for some $0<\lambda<1$. The second characteristic function of $X_{\mathrm{n}}$ is then $\gamma(\mathrm{t})=\sum_{i=1}^{n} \mathrm{a}_{\mathrm{in}} \mathrm{t} \cdot \zeta^{\prime}\left(\lambda_{\mathrm{in}} \mathrm{a}_{\mathrm{in}} \mathrm{t}\right)$, implying $|\gamma(\mathrm{t})| \leq \sum_{i=1}^{n}\left|\mathrm{a}_{\mathrm{in}} \mathrm{t} \cdot \zeta^{\prime}\left(\lambda_{\mathrm{in}} \mathrm{a}_{\mathrm{in}} \mathrm{t}\right)\right| \leq$ $|\mathrm{t}| \cdot\left(\max _{\mathrm{i} \leq \mathrm{n}}\left|\zeta^{\prime}\left(\lambda_{\mathrm{in}} \mathrm{a}_{\mathrm{in}} \mathrm{t}\right)\right|\right) \cdot \sum_{i=1}^{n}\left|\mathrm{a}_{\mathrm{in}}\right|$. Then $\lim \sum_{i=1}^{n}\left|\mathrm{a}_{\mathrm{in}}\right|<\infty$ and $\lim \left(\max _{\mathrm{i} \leq \mathrm{n}}\left|\mathrm{a}_{\mathrm{in}}\right|\right)=0$ imply $\gamma(\mathrm{t}) \longrightarrow 0$ for each t , and hence $\mathrm{X}_{\mathrm{n}}$ converges in distribution, hence in probability, to 0 .

Theorem 4.16. Assume random variables $\mathrm{Y}_{\mathrm{i}}$ are i.i.d. with mean zero and variance $\sigma^{2} \in(0,+\infty)$. If an array $\mathrm{a}_{\mathrm{in}}$ satisfies $\lim _{\mathrm{n} \rightarrow \infty} \max _{\mathrm{j} \leq \mathrm{n}}\left|\mathrm{a}_{\mathrm{jn}}\right|=0$ and $\lim _{\mathrm{n}-\infty} \sum_{i=1}^{n} \mathrm{a}_{\mathrm{in}}{ }^{2}=1$, then $\mathrm{X}_{\mathrm{n}}{ }_{\mathrm{d}} \mathrm{X}_{\mathrm{o}} \sim N\left(0, \sigma^{2}\right)$.

Proof: The argument parallels the Lindeberg-Levy CLT proof. The second characteristic function of $X_{n}$ has the Taylor's expansion $\gamma(\mathrm{t})=-(1 / 2) \sigma^{2} \mathrm{t}^{2} \mathrm{a}_{\mathrm{in}}+\left[\zeta^{\prime \prime}\left(\lambda_{\mathrm{in}} \mathrm{a}_{\mathrm{in}} \mathrm{t}\right)+\sigma^{2}\right] \cdot \mathrm{a}_{\mathrm{in}}{ }^{2} \mathrm{t}^{2} / 2$, where $\lambda_{\mathrm{in}} \in(0,1)$. The limit assumptions imply $\gamma(\mathrm{t})+(1 / 2) \sigma^{2} \mathrm{t}^{2}$ is bounded in magnitude by

$$
\sum_{i=1}^{n}\left|\zeta "\left(\lambda_{\mathrm{in}} \mathrm{a}_{\mathrm{in}} \mathrm{t}\right)+\sigma^{2}\right| \cdot \mathrm{a}_{\mathrm{in}} \mathrm{t}^{2} / 2 \leq\left[\begin{array}{ll}
\sum_{i=1}^{n} & \left.\mathrm{a}_{\mathrm{in}}{ }^{2} \mathrm{t}^{2} / 2\right] \cdot \max _{\mathrm{inn}}\left|\zeta\left(\zeta_{\mathrm{in}} \mathrm{a}_{\mathrm{in}} \mathrm{t} \mathrm{t}\right)+\sigma^{2}\right| .
\end{array}\right.
$$

This converges to zero for each $t$ since $\lim _{\mathrm{n} \rightarrow \infty} \max _{\mathrm{i} \leq \mathrm{n}}\left|\zeta^{\prime \prime}\left(\lambda_{\mathrm{in}} \mathrm{a}_{\mathrm{in}} \mathrm{t}\right)+\sigma^{2}\right| \rightarrow 0$. Therefore, $\gamma(\mathrm{t})$ converges to the characteristic function of a normal with mean 0 and variance $\sigma^{2}$.
4.5.3. The limit theorems 4.13 and 4.14 are special cases of a limit theory for what are called triangular arrays of random variables, $\mathrm{Y}_{\mathrm{nt}}$ with $\mathrm{t}=1,2, \ldots, \mathrm{n}$ and $\mathrm{n}=1,2,3, \ldots$. (One additional level of generality could be introduced by letting $t$ range from 1 up to a function of $n$ that increases to infinity, but this is not needed for most applications.) This setup will include simple cases like $\mathrm{Y}_{\mathrm{nt}}$ $=Z_{l} / n$ or $Y_{n t}=Z_{l} / n^{1 / 2}$, and more general weightings like $Y_{n t}=a_{n t} Z_{t}$ with an array of constants $a_{n t}$, but can also cover more complicated cases. We first give limit theorems for $Y_{n t}$ that are uncorrelated or independent within each row. These are by no means the strongest obtainable, but they have the merit of simple proofs.

Theorem 4.17. Assume random variables $Y_{n t}$ for $t=1,2, \ldots, n$ and $n=1,2,3, \ldots$ are uncorrelated across t for each n, with $\mathbf{E} \mathrm{Y}_{\mathrm{nt}}=0, \mathbf{E} \mathrm{Y}_{\mathrm{nt}}{ }^{2}=\sigma_{\mathrm{nt}}{ }^{2}$. Then, $\quad \sum_{i=1}^{n} \quad \sigma_{\mathrm{nt}}{ }^{2} \rightarrow 0$ implies $\quad \sum_{i=1}^{n} \quad \mathrm{Y}_{\mathrm{nt}} \overrightarrow{\mathrm{p}} 0$.

Proof: Apply Chebyshev's inequality.
Theorem 4.18. Assume random variables $Y_{n t}$ for $t=1,2, \ldots, n$ and $n=1,2,3, \ldots$ are independent across t for each n, with $\mathbf{E} \mathrm{Y}_{\mathrm{nt}}=0, \mathbf{E} \mathrm{Y}_{\mathrm{nt}}{ }^{2}=\sigma_{\mathrm{nt}}{ }^{2}, \quad \sum_{i=1}^{n} \sigma_{\mathrm{nt}}{ }^{2} \rightarrow 1, \quad \sum_{i=1}^{n} \mathbf{E}\left|\mathrm{Y}_{\mathrm{nt}}\right|^{3} \rightarrow 0$, and $\sum_{i=1}^{n} \quad \sigma_{\mathrm{nt}}{ }^{4} \rightarrow 0$. Then $\mathrm{X}_{\mathrm{n}}=\sum_{i=1}^{n} \quad \mathrm{Y}_{\mathrm{nt}}{ }_{\mathrm{d}} \mathrm{X}_{\mathrm{o}} \sim N(0,1)$.

Proof: From the properties of characteristic functions (see 3.5.12), the c.f. of $\mathrm{Y}_{\mathrm{nt}}$ has a Taylor's expansion that satisfies $\left|\psi_{n t}(\mathrm{~s})-1+\mathrm{s}^{2} \sigma_{\mathrm{nt}}{ }^{2} / 2\right| \leq|\mathrm{s}|^{3} \mathbf{E}\left|\mathrm{Y}_{\mathrm{nt}}\right|^{3 / 6}$. Therefore, the c.f. $\gamma_{\mathrm{n}}(\mathrm{s})$ of $X_{\mathrm{n}}$ satisfies $\log \gamma_{\mathrm{n}}(\mathrm{s})=\sum_{i=1}^{n} \log \left(1-\mathrm{s}^{2} \sigma_{\mathrm{nt}}^{2} / 2+\lambda_{\mathrm{nt}}|\mathrm{s}|^{3} \mathbf{E}\left|\mathrm{Y}_{\mathrm{nt}}\right|^{3} / 6\right)$, where $\left|\lambda_{\mathrm{nt}}\right| \leq 1$. From 2.1.10, we have the inequality that for $|a|<1 / 3$ and $|b|<1 / 3,|\log (1+a+b)-a|<4|b|+3|a|^{2}$. Then, the assumptions guarantee that $\left|\log \gamma_{\mathrm{n}}(\mathrm{s})+\mathrm{s}^{2} \sum_{i=1}^{n} \quad \sigma_{\mathrm{nt}}{ }^{2} / 2\right| \leq 4|\mathrm{~s}|^{3} \sum_{i=1}^{n} \quad \mathbf{E}\left|\mathrm{Y}_{\mathrm{nt}}\right|^{3} / 6+3 \mathrm{~s}^{4} \sum_{i=1}^{n} \sigma_{\mathrm{nt}}{ }^{4} / 4$. The assumptions then imply that $\log \gamma_{n}(s) \rightarrow-s^{2} / 2$, establishing the result.

In the last theorem, note that if $Y_{n t}=n^{-1 / 2} Z_{t}$, then $E\left|Z_{t}\right|^{3}$ bounded is sufficient to satisfy all the assumptions. Another set of limit theorems can be stated for triangular arrays with the property that the random variables within each row form a martingale difference sequence. Formally, consider random variables $\mathrm{Y}_{\mathrm{nt}}$ for $\mathrm{t}=1, \ldots, \mathrm{n}$ and $\mathrm{n}=1,2,3, \ldots$ that are adapted to $\sigma$-fields $\boldsymbol{G}_{\mathrm{nt}}$ that are a filtration in t for each n , with the property that $\mathrm{E}\left\{\mathrm{Y}_{\mathrm{n} t} \mid \boldsymbol{G}_{\mathrm{n}, \mathrm{t}-1}\right\}=0$; this is called a martingale difference array. A WLLN for this case is adapted from J. Davidson (1994), p. 299.

Theorem 4.19. If $\mathrm{Y}_{\mathrm{nt}}$ and $\boldsymbol{G}_{\mathrm{nt}}$ for $\mathrm{t}=1, \ldots, \mathrm{n}$ and $\mathrm{n}=1,2,3, \ldots$ is an adapted martingale difference array with $\left|\mathrm{Y}_{\mathrm{nt}}\right| \leq \mathrm{M}, \mathbf{E} \mathrm{Y}_{\mathrm{nt}}{ }^{2}=\sigma_{\mathrm{nt}}{ }^{2}, \quad \sum_{i=1}^{n} \quad \sigma_{\mathrm{nt}}$ uniformly bounded, and $\sum_{i=1}^{n} \sigma_{\mathrm{nt}}{ }^{2} \rightarrow 0$, then

$$
\sum_{i=1}^{n} \quad \mathrm{Y}_{\mathrm{nt}} \overrightarrow{\mathrm{p}}_{\mathrm{p}} 0
$$

The following CLT for martingale difference arrays is taken from D. Pollard (1984), p. 170-174.
Theorem 4.20. If $\mathrm{Y}_{\mathrm{nt}}$ and $\boldsymbol{G}_{\mathrm{nt}}$ for $\mathrm{t}=1, \ldots, \mathrm{n}$ and $\mathrm{n}=1,2,3, \ldots$ is an adapted martingale difference array, $\lambda_{\mathrm{nt}}{ }^{2}=\mathrm{E}\left(\mathrm{Y}_{\mathrm{nt}}{ }^{2} \mid \boldsymbol{G}_{\mathrm{n}, \mathrm{t}-1}\right)$ is the conditional variance of $\mathrm{Y}_{\mathrm{nt}}, \quad \sum_{i=1}^{n} \quad \lambda_{\mathrm{nt}}{ }^{2} \rightarrow_{\mathrm{p}} \sigma^{2} \in(0,+\infty)$, and if for each $\varepsilon>0, \quad \sum_{t=1}^{n} \quad \mathbf{E} \mathrm{Y}_{\mathrm{nt}}{ }^{2} \cdot \mathbf{1}\left(\left|\mathrm{Y}_{\mathrm{nt}}\right|>\varepsilon\right) \rightarrow 0$, then $\mathrm{X}_{\mathrm{n}}=\sum_{i=1}^{n} \quad \mathrm{Y}_{\mathrm{nt}}{ }_{\mathrm{d}} \mathrm{X}_{\mathrm{o}} \sim N\left(0, \sigma^{2}\right)$.
4.5.4. Chebyshev's inequality gives an easy, but crude, bound on the probability in the tail of a density. For random variables with well behaved tails, there are sharper bounds that can be used to get sharper limit theorems. The following inequality, due to Hoeffding, is one of a series of results called exponential inequalities that are stated and proved in D. Pollard (1984), p. 191-193: If $Y_{n}$ are independent random variables with zero means that satisfy the bounds $-a_{n} \leq Y_{n} \leq b_{n}$, then $\mathrm{P}\left(n^{-1} \sum_{i=1}^{n} \quad \mathrm{Y}_{\mathrm{i}} \geq \varepsilon\right) \leq \exp \left(-2 \mathrm{n}^{2} \varepsilon^{2} / \sum_{i=1}^{n}\left(\mathrm{~b}_{\mathrm{i}}+\mathrm{a}_{\mathrm{i}}\right)^{2}\right)$. Note that in Hoeffding's inequality, if $\left|\mathrm{Y}_{\mathrm{n}}\right| \leq$ M , then $\mathrm{P}\left(\left|n^{-1} \sum_{i=1}^{n} \quad \mathrm{Y}_{\mathrm{i}}\right| \geq \varepsilon\right) \leq 2 \cdot \exp \left(-\mathrm{n} \varepsilon^{2} / 2 \mathrm{M}^{2}\right)$. The next theorem gets a strong law of large numbers with weaker than usual scaling:

Theorem 4.21. If $Y_{n}$ are independent random variables with zero means and $\left|Y_{n}\right| \leq M$, then $X_{n}$ $=n^{-1} \sum_{i=1}^{n} \quad \mathrm{Y}_{\mathrm{i}}$ satisfies $\mathrm{X}_{\mathrm{k}} \cdot \mathrm{k}^{1 / 2} / \log (\mathrm{k}) \rightarrow_{\text {as }} 0$.

Proof: Hoeffding's inequality implies $\operatorname{Prob}\left(\mathrm{k}^{1 / 2}\left|\mathrm{X}_{\mathrm{k}}\right|>\varepsilon \cdot \log \mathrm{k}\right)<2 \cdot \exp \left(-(\log \mathrm{k}) \varepsilon^{2} / 2 \mathrm{M}^{2}\right)$, and hence

$$
\begin{aligned}
& \sum_{k=n+1}^{\infty} \operatorname{Prob}\left(\mathrm{k}^{1 / 2}\left|\mathrm{X}_{\mathrm{k}}\right|>\varepsilon \cdot \log \mathrm{k}\right) \leq \int_{z=n}^{\infty} 2 \cdot \exp \left(-(\log \mathrm{z})^{2} \varepsilon^{2} / 2 \mathrm{M}^{2}\right) \mathrm{dz} \\
& \leq(6 / \varepsilon) \cdot \exp \left(\mathrm{M}^{2} / 2 \varepsilon^{2}\right) \cdot \Phi(-\varepsilon \cdot(\log \mathrm{n}) / \mathrm{M}+\mathrm{M} / \varepsilon)
\end{aligned}
$$

with the standard normal CDF $\Phi$ resulting from direct integration. Applying Theorem 4.2, this inequality implies $n^{1 / 2}\left|X_{n}\right| / \log n \rightarrow$ as 0 .

If the $\mathrm{Y}_{\mathrm{i}}$ are not necessarily bounded, but have a proper moment generating function, one can get an exponential bound from the moment generating function.

Theorem 4.22. If i.i.d. mean-zero random variables $Y_{i}$ have a proper moment generating function, then $X_{n}=n^{-1} \sum_{i=1}^{n} \quad Y_{i}$ satisfies $\mathrm{P}\left(\mathrm{X}_{\mathrm{n}}>\varepsilon\right)<\exp \left(-\tau \varepsilon \mathrm{n}^{1 / 2}+\kappa\right)$, where $\tau$ and $\kappa$ are positive constants determined by the distribution of $\mathrm{Y}_{\mathrm{i}}$.

Proof: $\mathrm{P}(\mathrm{Z}>\varepsilon)=\int_{z>\varepsilon} \mathrm{F}(\mathrm{dz}) \leq \int_{z>\varepsilon} \mathrm{e}^{(z-\varepsilon) \mathrm{F}} \mathrm{F}(\mathrm{dz}) \leq \mathrm{e}^{-\varepsilon \varepsilon} \mathrm{Ee}^{\mathrm{Zt}}$ for a random variable Z . Let $\mathrm{m}(\mathrm{t})$ be the moment generating function of $\mathrm{Y}_{\mathrm{i}}$ and $\tau$ be a constant such that $\mathrm{m}(\mathrm{t})$ is finite for $|\mathrm{t}|<2 \tau$. Then one has $m(t)=1+m^{\prime \prime}(\lambda t) t^{2} / 2$ for some $|\lambda|<1$, for each $|t|<2 \tau$, from the properties of $m g f$ (see 3.5.12).

The mgf of $X_{n}$ is $m(t / n)^{n}=\left(1+m^{"}(\lambda t / n) t^{2} / 2 n^{2}\right)^{n}$, finite for $|t| / n \leq 2 \tau$. Replace $t / n$ by $\tau n^{-1 / 2}$ and observe that $\mathrm{m}^{\prime \prime}(\lambda \mathrm{t} / \mathrm{n}) \leq \mathrm{m}^{\prime \prime}\left(\tau \mathrm{n}^{-1 / 2}\right)$ and $\left(1+\mathrm{m}^{\prime \prime}\left(\tau \mathrm{n}^{-1 / 2}\right) \tau^{2} / 2 \mathrm{n}\right)^{\mathrm{n}} \leq \exp \left(\mathrm{m}^{\prime \prime}\left(\tau \mathrm{n}^{-1 / 2}\right) \tau^{2} / 2\right)$. Substituting these expressions in the initial inequality gives $\mathrm{P}\left(\mathrm{X}_{\mathrm{n}}>\varepsilon\right) \leq \exp \left(-\tau \varepsilon \mathrm{n}^{1 / 2}+\mathrm{m}^{\prime \prime}\left(\tau \mathrm{n}^{-1 / 2}\right) \tau^{2} / 2\right)$, and the result holds with $\kappa=\mathrm{m}^{\prime \prime}(\tau) \tau^{2} / 2$.

Using the same argument as in the proof of Theorem 4.19 and the inequality $\mathrm{P}\left(\mathrm{X}_{\mathrm{n}}>\varepsilon\right)<$ $\exp \left(-\tau \varepsilon n^{1 / 2}+\kappa\right)$ from Theorem 4.20, one can show that $X_{k} \cdot \mathrm{k}^{1 / 2} /(\log k)^{2} \rightarrow_{\text {as }} 0$, a SLLN with weak scaling.
4.5.5. This section states a uniform SLLN for random functions on compact set $\Theta$ in a Euclidean space $\mathbb{R}^{k}$. Let $(\mathbf{S}, \boldsymbol{F}, \mathrm{P})$ denote a probability space. Define a random function as a mapping Y from $\Theta \times \mathrm{S}$ into $\mathbb{R}$ with the property that for each $\theta \in \Theta, \mathrm{Y}(\theta, \cdot)$ is measurable with respect to ( $\mathrm{S}, F, \mathrm{P}$ ). Note that $\mathrm{Y}(\theta, \cdot)$ is simply a random variable, and that $\mathrm{Y}(\cdot, \mathrm{s})$ is simply a function of $\theta \in \Theta$. Usually, the dependence of Y on the state of nature is suppressed, and we simply write $\mathrm{Y}(\theta)$. A random function is also called a stochastic process, and $\mathrm{Y}(\cdot, \mathrm{s})$ is termed a realization of this process. A random function $\mathrm{Y}(\theta, \cdot)$ is almost surely continuous at $\theta_{\mathrm{o}} \in \Theta$ if for s in a set that occurs with probability one, $\mathrm{Y}(\cdot, \mathrm{s})$ is continuous in $\theta$ at $\theta_{0}$. It is useful to spell out this definition in more detail. For each $\varepsilon>0$, define $\mathbf{A}_{\mathrm{k}}\left(\varepsilon, \theta_{o}\right)=\left\{s \in \boldsymbol{S}\left|\underset{\left|\theta-\theta_{o}\right| \leq 1 / k}{\sup }\right| Y(\theta, s)-Y\left(\theta_{o}, s\right) \mid>\varepsilon\right\}$. Almost sure continuity states that these sets converge monotonically as $\mathrm{k} \rightarrow \infty$ to a set $\mathbf{A}_{\mathrm{o}}\left(\varepsilon, \theta_{\mathrm{o}}\right)$ that has probability zero.

The condition of almost sure continuity allows the modulus of continuity to vary with s , so there is not necessarily a fixed neighborhood of $\theta_{0}$ independent of $s$ on which the function varies by less than $\varepsilon$. For example, the function $\mathrm{Y}(\theta, \mathrm{s})=\theta^{s}$ for $\theta \in[0,1]$ and s uniform on $[0,1]$ is continuous at $\theta=0$ for every s, but $\mathbf{A}_{\mathrm{k}}(\varepsilon, 0)=[0,(-\log \varepsilon) /(\log \mathrm{k}))$ has positive probability for all k . The exceptional sets $\mathbf{A}_{\mathrm{k}}(\varepsilon, \theta)$ can vary with $\theta$, and there is no requirement that there be a set of s with probability one, or for that matter with positive probability, where $\mathrm{Y}(\theta, \mathrm{s})$ is continuous for all $\theta$. For example, assuming $\theta \in[0,1]$ and $s$ uniform on $[0,1]$, and defining $Y(\theta, s)=1$ if $\theta \geq s$ and $Y(\theta, s)=0$ otherwise gives a function that is almost surely continuous everywhere and always has a discontinuity.

Theorem 4.3 in Section 4.1 established that convergence in probability is preserved by continuous mappings. The next result extends this to almost surely continuous transformations; the result below is taken from Pollard (1984), p. 70.

Theorem 4.23. (Continuous Mapping). If $Y_{n}(\theta) \rightarrow_{p} Y_{o}(\theta)$ uniformly for $\theta$ in $\Theta \subseteq \mathbb{R}^{k}$, random vectors $\tau_{0}, \tau_{n} \in \Theta$ satisfy $\tau_{\mathrm{n}} \overrightarrow{\mathrm{p}} \tau_{\mathrm{o}}$, and $\mathrm{Y}_{\mathrm{o}}(\theta)$ is almost surely continuous at $\tau_{\mathrm{o}}$, then $\mathrm{Y}_{\mathrm{n}}\left(\tau_{\mathrm{n}}\right) \overrightarrow{\mathrm{p}}_{\mathrm{p}} \mathrm{Y}_{\mathrm{o}}\left(\tau_{\mathrm{o}}\right)$.

Consider i.i.d. random functions $\mathrm{Y}_{\mathrm{i}}(\theta)$ that have a finite mean $\psi(\theta)$ for each $\theta$, and consider the average $\mathrm{X}_{\mathrm{n}}(\theta)=n^{-1} \sum_{i=1}^{n} \quad \mathrm{Y}_{\mathrm{i}}(\theta)$. Kolmogorov's SLLN I implies that pointwise, $\mathrm{X}_{\mathrm{n}}(\theta) \rightarrow_{\mathrm{as}} \psi(\theta)$.

However, we sometimes need in statistics a stronger result that $X_{n}(\theta)$ is uniformly close to $\psi(\theta)$ over the whole domain $\Theta$. This is not guaranteed by pointwise convergence. For example, the random function $\mathrm{Y}_{\mathrm{n}}(\mathrm{s}, \theta)=1$ if $\mathrm{n}^{2} \cdot|\mathrm{~s}-\theta| \leq 1$, and $\mathrm{Y}_{\mathrm{n}}(\mathrm{s}, \theta)=0$ otherwise, where the sample space is the unit interval with uniform probability, has $\mathrm{P}\left(\mathrm{Y}_{\mathrm{n}}(\cdot, \theta)>0\right) \leq 2 / \mathrm{n}^{2}$ for each $\theta$. This is sufficient to give $\mathrm{Y}_{\mathrm{n}}(\cdot, \theta) \rightarrow_{\text {as }} 0$ pointwise. However, $\mathrm{P}\left(\sup _{\theta} \mathrm{Y}_{\mathrm{n}}(\theta)>0\right)=1$.

Theorem 4.24. (Uniform SLLN). Assume $\mathrm{Y}_{\mathrm{i}}(\theta)$ are independent identically distributed random functions with a finite mean $\psi(\theta)$ for $\theta$ in a closed bounded set $\Theta \subseteq \mathbb{R}^{\mathrm{k}}$. Assume $\mathrm{Y}_{\mathrm{i}}(\cdot)$ is almost surely continuous at each $\theta \in \Theta$. Assume that $Y_{i}(\cdot)$ is dominated; i.e., there exists a random variable $Z$ with a finite mean that satisfies $Z \geq \sup _{\theta \in \Theta}\left|Y_{1}(\theta)\right|$. Then $\psi(\theta)$ is continuous in $\theta$ and $\mathrm{X}_{\mathrm{n}}(\theta)=\frac{1}{n} \sum_{i=1}^{n} Y_{i}(\theta) \quad$ satisfies $\sup _{\theta \in \Theta}\left|\mathrm{X}_{\mathrm{n}}(\theta)-\psi(\theta)\right| \rightarrow_{\text {as }} 0$.

Proof: We follow an argument of Tauchen (1985). Let (S,F,P) be the probability space, and write the random function $\mathrm{Y}_{\mathrm{i}}(\theta, \mathrm{s})$ to make its dependence on the state of Nature explicit. We have $\psi(\theta)$
$=\int_{S} \mathrm{Y}(\theta, \mathrm{s}) \mathrm{P}(\mathrm{ds})$. Define $\mathrm{u}\left(\theta_{0}, \mathrm{~s}, \mathrm{k}\right)=\sup _{\left|\theta-\theta_{0}\right| \leq 1 / k}\left|\mathrm{Y}(\theta, \mathrm{s})-\mathrm{Y}\left(\theta_{0}, \mathrm{~s}\right)\right|$. Let $\varepsilon>0$ be given. Let
$\mathbf{A}_{\mathrm{k}}\left(\varepsilon / 2,0_{\mathrm{o}}\right)$ be the measurable set given in the definition of almost sure continuity, and note that for $\mathrm{k}=\mathrm{k}\left(\varepsilon / 2, \theta_{\mathrm{o}}\right)$ sufficiently large, the probability of $\mathbf{A}_{\mathrm{k}}\left(\varepsilon / 2, \theta_{\mathrm{o}}\right)$ is less than $\varepsilon /(4 \cdot \mathbf{E} Z)$. Then,

$$
\begin{gathered}
\mathbf{E u}\left(\theta_{o} \cdot \cdot, \mathrm{k}\right) \leq \int_{A_{k}\left(\varepsilon / 2, \theta_{o}\right)} \mathrm{u}\left(\theta_{\mathrm{o}}, \mathrm{~s}, \mathrm{k}\right) \mathrm{P}(\mathrm{ds})+\int_{A_{k}\left(\varepsilon / 2, \theta_{o}\right)^{c}} \mathrm{u}\left(\theta_{\mathrm{o}}, \mathrm{~s}, \mathrm{k}\right) \mathrm{P}(\mathrm{ds}) \\
\quad \leq \int_{A_{k}\left(\varepsilon / 2, \theta_{o}\right)} 2 \cdot \mathrm{Z}(\mathrm{~s}) \cdot \mathrm{P}(\mathrm{ds})+\int_{A_{k}\left(\varepsilon / 2, \theta_{o}\right)^{c}}(\varepsilon / 2) \cdot \mathrm{P}(\mathrm{ds}) \leq \varepsilon .
\end{gathered}
$$

Let $\mathbf{B}\left(\theta_{\mathrm{o}}\right)$ be an open ball of radius $1 / k\left(\varepsilon / 2, \theta_{\mathrm{o}}\right)$ about $\theta_{0}$. These balls constructed for each $\theta_{\mathrm{o}} \in \Theta$ cover the compact set $\Theta$, and it is therefore possible to extract a finite subcovering of balls $\mathbf{B}\left(\theta_{\mathrm{j}}\right)$ with centers at points $\theta_{\mathrm{j}}$ for $\mathrm{j}=1, \ldots, \mathrm{~J}$. Let $\mu_{\mathrm{j}}=\mathbf{E u}\left(\theta_{\mathrm{j}} ;, \mathrm{k}\left(\boldsymbol{\varepsilon} / 2, \theta_{\mathrm{j}}\right)\right) \leq \boldsymbol{\varepsilon}$. For $\theta \in \mathbf{B}\left(\theta_{\mathrm{j}}\right),\left|\psi(\theta)-\psi\left(\theta_{\mathrm{j}}\right)\right| \leq \mu_{\mathrm{j}} \leq$ $\varepsilon$. Then

$$
\begin{aligned}
\sup _{\theta \in B\left(\theta_{\mathrm{j}}\right)}\left|\mathrm{X}_{\mathrm{n}}(\theta)-\psi(\theta)\right| & \leq\left|\mathrm{X}_{\mathrm{n}}(\theta)-\mathrm{X}_{\mathrm{n}}\left(\theta_{\mathrm{j}}\right)-\mu_{\mathrm{j}}\right|+\mu_{\mathrm{j}}+\left|\mathrm{X}_{\mathrm{n}}\left(\theta_{\mathrm{j}}\right)-\psi\left(\theta_{\mathrm{j}}\right)\right|+\left|\psi\left(\theta_{\mathrm{j}}\right)-\psi(\theta)\right| \\
& \leq\left|\frac{1}{n} \sum_{i=1}^{n} \mathrm{u}\left(\theta_{\mathrm{j}} ;, \mathrm{k}\left(\varepsilon / 2, \theta_{\mathrm{j}}\right)\right)-\mu_{\mathrm{j}}\right|+\varepsilon+\left|\mathrm{X}_{\mathrm{n}}\left(\theta_{\mathrm{j}}\right)-\psi\left(\theta_{\mathrm{j}}\right)\right|+\varepsilon
\end{aligned}
$$

Apply Kolmogorov's SLLN to each of the first and third terms to determine a sample size $\mathrm{n}_{\mathrm{j}}$ such that

$$
\mathrm{P}\left(\sup _{n \geq n_{j}}\left|n^{-1} \sum_{i=1}^{n} \mathrm{u}\left(\theta_{\mathrm{j}} ;, \mathrm{k}\left(\varepsilon / 2, \theta_{\mathrm{j}}\right)\right)-\mu_{\mathrm{j}}\right|>\varepsilon\right)<\varepsilon / 2 \mathrm{~J}
$$

and

$$
\mathrm{P}\left(\sup _{n \geq n_{j}}\left|\mathrm{X}_{\mathrm{n}}\left(\theta_{\mathrm{j}}\right)-\psi\left(\theta_{\mathrm{j}}\right)\right|>\varepsilon\right)<\varepsilon / 2 \mathrm{~J} .
$$

With probability at least $1-\varepsilon / \mathrm{J}, \quad \sup _{\theta \in B\left(\mathrm{j}_{j}\right)}\left|\mathrm{X}_{\mathrm{n}}(\theta)-\psi(\theta)\right| \leq 4 \varepsilon$. Then, with probability at least $1-\varepsilon$,

$$
\sup _{\theta \in \Theta}\left|X_{\mathrm{n}}(\theta)-\psi(\theta)\right| \leq 4 \varepsilon \text { for } \mathrm{n}>\mathrm{n}_{\mathrm{o}}=\max \left(\mathrm{n}_{\mathrm{j}}\right) .
$$

The construction in the proof of the theorem of a finite number of approximating points can be reinterpreted as the construction of a finite family of functions, the $\mathrm{Y}\left(\theta_{\mathrm{j}}, \cdot\right)$, with the approximation property that the expectation of the absolute difference between $\mathrm{Y}(\theta, \cdot)$ for any $\theta$ and one of the members of this finite family is less than $\varepsilon$. Generalizations of the uniform SLLN above can be obtained by recognizing that it is this approximation property that is critical, with a limit on how rapidly the size of the approximating family can grow with sample size for a given $\varepsilon$, rather than continuity per se; see D. Pollard (1984).

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