

Stochastonomicon

Jyotirmoy Bhattacharya

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Preface

You should not read this. Here's why:

1. This is work in progress and is likely to contain many errors.
2. There is nothing original here. Everything in these notes I found in a book or on the Internet and at best rearranged.
3. Thinking through stuff on your own is likely to be more productive.

But if you want to see a poor fellow trudge his way painfully through the landscape of modern probability, welcome!

You can download the latest version of this document from <https://www.jyotirmoy.net/misc/probability-guide.pdf>.

Please send comments and corrections to jyotirmoy@jyotirmoy.net.

Chapter 1

Weak Convergence

1.1 Weak Convergence

If we say that a sequence of measures μ_n converges to a measure μ , what should we require? The following come to mind:

1. The measure of each measurable set should converge uniformly.
2. The measure of each measurable set should converge.
3. (If the state space is a topological space) The expectation of each continuous function of the state should converge.
4. (If the state space is a topological space) The expectation of each continuous bounded function of the state should converge.

In the classical central limit theorem setting we can only prove weak convergence [TODO: Give example]

In the Markov chain setting we can do better and provide TV convergence (give citation).

For economic problems weak convergence should be enough? [TODO: What does this even mean?]

1.2 Unbounded functions

If $\mu_n \Rightarrow \mu$ then can we say that even for unbounded functions g we must have

$$\int g d\mu_n \rightarrow \int g d\mu$$

No. Take $\mu_n = (1 - 1/n)\delta_0 + (1/n)\delta_n$ and $g(x) = x$.

Chapter 2

The Multivariate Normal Distribution

Jacod and Protter's *Probability Essentials* has a nice chapter on the multivariate normal distribution.

2.1 Conditional distribution

Let (x, y) be a \mathbb{R}^{n+m} dimensional multivariate normal vector where x is n -dimensional and y is m -dimensional respectively. Let the mean of (x, y) be $(0, 0)$ and its variance-covariance matrix be nonsingular.

What is the conditional distribution of x given y ?

The answer is easy if x and y are independent, because then the conditional distribution of x is just the unconditional distribution, which is a multivariate normal distribution because $a'x = (a, 0)'(x, y)$ is normal for every choice of a .

What if x and y are not independent. Our approach will be to reduce to the easier independent case. We know that for components of a multivariate normal vector, zero covariance implies independence. We use the Gram-Schmidt trick of subtracting a suitably scaled copy of one variable from another to get zero covariance.

Define

$$w = y - Ax$$

We want $E[wx'] = 0$ which gives us

$$A = E(yx')[E(xx')]^{-1}$$

We know that $E(xx')$ is nonsingular since if it were not there would be a non-zero a such that $E(xx')a = 0$, so that $E[(a'x)(a'x)'] = 0$, so that $(a, 0)'E[(x, y)(x, y)'](a, 0) = 0$, which is inconsistent with $E[(x, y)(x, y)']$ being nonsingular as assumed.

Since (x, w) is obtained as a linear transformation of (x, y) the former is also a multivariate normal vector. Moreover, x and w have zero covariance

by construction and hence are independent. So the conditional distribution of w given x is just the unconditional distribution of w which is just a normal.

Now $y = w + Ax$, so conditional on x , y is just w shifted by a constant, so if w has a conditional normal distribution, so has y .

To pin down the exact distribution of w we compute its mean and variance-covariance matrix.

For the mean, since w and x are independent

$$\begin{aligned}\mu_{y|x} &= E(y | x) \\ &= E(w | x) + Ax\end{aligned}$$

since w and x are independent,

$$= E(w) + Ax$$

since $E(w) = E(y) = 0$

$$\begin{aligned}&= 0 + Ax \\ &= E(yx')[E(xx')]^{-1}x\end{aligned}$$

For the conditional variance covariance

$$\begin{aligned}\Sigma_{y|x} &= E\{[y - E(y | x)]'[y - E(y | x)] | x\} \\ &= E\{(y - Ax)'(y - Ax) | x\} \\ &= E(ww' | x)\end{aligned}$$

since w is independent of x

$$\begin{aligned}&= E[ww'] \\ &= E[(y - Ax)(y - Ax)'] \\ &= E(yy') - AE(xy') - E(yx')A + AE(xx')A\end{aligned}$$

substituting the value of A

$$= E(yy') - E(yx')[E(xx')]^{-1}E(xy')$$

Appendix A

Metric Spaces

A.1 The Distance Function

Where not otherwise specified X is a metric space with metric $d(\cdot, \cdot)$.

Definition A.1. If F is a non-empty closed set and $x \in X$ we define

$$d(x, F) = \inf_{y \in F} d(x, y)$$

Proposition A.2. For a fixed F , $d(x, F)$ is an uniformly continuous function of F . Indeed,

$$|d(x, F) - d(y, F)| \leq d(x, y)$$

Proposition A.3. If F is a non-empty closed set and $x \notin F$ then

$$d(x, F) > 0$$

Definition A.4. If F_1 and F_2 are non-empty disjoint closed sets, we define

$$d(F_1, F_2) = \inf_{x \in F_1} d(x, F_2)$$

Proposition A.5. Under conditions of Definition A.4, $d(F_1, F_2) > 0$.

A.2 Urysohn's Theorem

Theorem A.6 (Urysohn). If F_1 and F_2 are disjoint closed sets then there exists a continuous function f such that $f(x) = 0$ for $x \in F_1$ and $f(x) = 1$ for $x \in F_2$.

Proof. The following function has these properties

$$f(x) = \frac{d(x, F_1)}{d(x, F_1) + d(x, F_2)}.$$

□

Proposition A.7. *If F is a closed set and U is an open set such that $F \subset U$ then it is possible to find an open set V such that*

$$F \subset V \subset \bar{V} \subset U$$

Proof. If F is empty we can take $V = \emptyset$. If $U = X$ we can take $V = X$.

Otherwise, F and U^c are non-empty, disjoint closed sets. Let $f(x)$ be the continuous function provided by Theorem A.6 which is 0 on F and 1 on U^c and take

$$V = \{x \mid f(x) < 1/2\}.$$

Then $F \subset V$ and by the continuity of f , V is an open set. $\bar{V} \subset \{x \mid f(x) \leq 1/2\}$ and the latter is disjoint from U^c so that $\bar{V} \subset U$. \square

A.3 Partitions of Unity

Theorem A.8. *Let U_i , $i = 1, \dots, N$ be an open cover of X . We can find continuous functions f_i , $i = 1, \dots, N$ such that $0 \leq f_i \leq 1$, $\text{supp } f_i \subset U_i$ and*

$$\sum_i f_i(x) = 1, \quad \text{for all } x \in X$$

In this case $\{f_i\}$ is called a partition of unity subordinate to $\{U_i\}$.

If we did not want $\text{supp } f_i \subset U_i$ and could instead have done just with $f_i(x) = 0$ for $x \notin U_i$, it would have been enough to take

$$f_i(x) = \frac{d(x, U_i^c)}{\sum_{j=1}^N d(x, U_j^c)}.$$

But the stronger conclusion in the theorem as stated above needs some more work. The following argument is taken from [lecture notes](#) by Marius Crainic.

We first prove the following following

Proposition A.9 (Refinement). *If $\{U_i\}$ is a finite open cover of X , then there is another open cover $\{V_i\}$ of X such that $\bar{V}_i \subset U_i$ for all i .*

Proof. Consider $F = X \setminus \bigcup_{j \neq 1} U_j$. Then F is closed, and it is a subset of U_1 (since the U_i cover X). By Proposition A.7 we can find an open V_1 such that $F \subset V_1 \subset \bar{V}_1 \subset U_1$.

From the definition of F and the fact that $F \subset V_1$ it follows that $\{V_1, U_2, \dots, U_n\}$ is also an open cover of X . We now repeat the argument in the para above with this new open cover to replace U_2 with V_2 and so on. \square

We are now in a position to prove the main theorem.

Proof of Theorem A.8. Let $\{U_i\}$ be the original open cover. We apply Proposition A.9 to create a refinement $\{V_i\}$ and then a further refinement $\{W_i\}$ such that $\bar{W}_i \subset V_i \subset \bar{V}_i \subset U_i$. Taking \bar{W}_i and V_i^c as the two sets in Theorem A.6 we can find a continuous function g_i such that $g_i(x) = 1$ on \bar{W}_i and $g_i(x) = 0$ on V_i^c . It follows that $\text{supp } g_i \subset \bar{V}_i \subset U_i$. Now $\sum_{i=1}^N g_i(x) \neq 0$ for all x since W_i is a cover of X and $g_i(x) = 1$ on W_i . So defining

$$f_i(x) = \frac{g_i(x)}{\sum_{j=1}^N g_j(x)}$$

fulfils all the requirements of the theorem. \square

Note: For a compact K , the standard topological theorem for LCH spaces gives a partition of unity where the functions have compact supports.

A.4 The space $C(X)$

We use $C(X)$ to denote the space of *bounded*, real-valued, continuous functions on X .

Theorem A.10. *The space $C(X)$ is separable if and only if X is compact.*

If. Adapted from Conway's *Functional Analysis*.

Suppose X is compact. For every n , X can be covered by a finite set of balls of radius $1/n$. Denote these by U_n^k for $k = 1, \dots, K_n$. Let $\{f_n^k\}$ be a partition of unity subordinate to $\{U_n^k\}$ for a given n , which exists by Theorem A.8. Now consider $\bigcup_n \{f_n^k \mid k = 1, \dots, K_n\}$. This is a countable set. Let \mathcal{F} be the set of rational finite linear combinations of functions from the set $\bigcup_n \{f_n^k \mid k = 1, \dots, K_n\}$. Then \mathcal{F} is a countable set of continuous functions. We will show that it is dense in $C(X)$.

Consider some $\phi(x) \in C(X)$. Since X is compact, f is uniformly continuous. Take $\epsilon > 0$ as given. We can find a n such that $d(x, y) < 2/n$ implies $|\phi(x) - \phi(y)| < \epsilon$.

Pick an arbitrary x_k in U_n^k for every $k = 1, \dots, K_n$ for the choice of n above. For each x_k pick a rational number q_k such that $|q_k - \phi(x_k)| < \epsilon$. Define

$$g(x) = \sum_{k=1}^{K_n} q_k f_n^k(x)$$

which is an element of \mathcal{F} .

Since the f_n^k sum up to 1, we can write

$$\phi(x) = \sum_{k=1}^{K_n} \phi(x) f_n^k(x),$$

so that we have

$$|\phi(x) - g(x)| \leq \sum_{k=1}^{K_n} |q_k - \phi(x)| f_n^k(x)$$

Now $f_n^k(x)$ is nonzero only on U_n^k and here

$$|q_k - \phi(x)| \leq |q_k - \phi(x_k)| + |\phi(x_k) - \phi(x)| \leq 2\epsilon$$

from which we conclude, again making use of the fact that f_n^k sum to 1, that

$$|\phi(x) - g(x)| \leq 2\epsilon.$$

Since the bound is independent of x our theorem is proved. □

An alternative proof uses the fact that X compact means that it is separable. Let $\{x_1, x_2, \dots\}$ be a countably dense subset of X . Define $h_n(x) = d(y, x_n)$ and consider the algebra of real polynomials in the $h_n(x)$. Apply the Stone-Weierstrass theorem to show that this algebra is dense in $C(X)$. Then approximate this algebra by the algebra of rational polynomials in $h_n(x)$ which is countable.

Only If. Again from Conway.

If X is not compact then there is a countably infinite subset $D = \{x_1, x_2, \dots\}$ such that every subset of D is closed. For any subset A of D we can use Theorem A.6 to find an $f_A(x) \in C(X)$ which is 1 on A and 0 on $D \setminus A$. For $A \neq B$, the distance between f_A and f_B in the uniform metric is at least 1. Since there are uncountably many subsets of D there are uncountably many elements in $C(X)$ such that the distance between any pair of them is at least 1. So no countable subset of $C(X)$ can be dense. □