Introduction

My notes from the legendary book of Papoulis and Pillai [1]. The book assumes that the reader does already have fairly strong background in calculus and linear algebra. I used two books [2, 3] as assistant for evaluating some of the mathematical identities or integrals.

This book made me realize that a good technical book is not necessarily one that's easy to follow – one that shows you very easily how each identity or result is arrived at. Some of the results in the book are presented without much explanation, and this required me to show a lot of effort to understand them; this effort seems to lead to a staying power.

Chapter 1

Sequences of Random Variables

1.1 Notes

- Multivariate Transformation
- How to integrate some RVs from a multi-variate distro to obtain ... (Sec 7.2)
- How to compute the mean conditioned on a subset of RVs
- Characteristic function leads to so much interesting applications, such as:
- Computing the PDF of Bernoulli from Binomials (#256)
- Computing the PDF of Poisson from Binomials (#256)
- The sum of jointly normal RVs is normal (#257)
- The sum of the squares of independent normal variables is chi-square (#259)
- The sum of two chi-square distros is also chi-square (#260)
- The optimal single-value estimation (in the MS sense), c, of a future value of a RV y is c = E(y).
- The optimal functional estimation of \mathbf{y} i.t.o. a (dependent) RV \mathbf{x} is $c(x) = E(\mathbf{y}|x)$. In general, this estimation is non-linear.
- The optimal linear f. estimation of y i.t.o. x is c(x) = Ax + B where $B = \eta_y A\eta_x$ and $A = r\sigma_x/\sigma_y$.
- For Gaussian RVs, linear and non-linear MS estimators are identical (#264).
- The orthogonality principle: The error between of linear estimator $\hat{\mathbf{y}} = A\mathbf{x} + B$ of \mathbf{y} is orthogonal to the data: $E\{(\mathbf{y} \hat{\mathbf{y}})\mathbf{x}\} = 0$
- Generalizes to the linear estimate of \mathbf{s} i.t.o. multi RVs $\mathbf{x}_1, \dots, \mathbf{x}_n$: $E\{(\mathbf{s} \hat{\mathbf{s}})\mathbf{x}_i\}$ for $i = 1, \dots, n$.
- Generalizes to non-linear estimation: $E\{[\mathbf{s} g(\mathbf{X})]w(\mathbf{X})\}$, where $g(\mathbf{X})$ is the non-linear MS estimator and $w(\mathbf{X})$ any function of data $\mathbf{X} = [\mathbf{x}_1, \dots, \mathbf{x}_n]$. (#269)
- Computing the linear MS estimator is much easier if the RVs \mathbf{x}_i are orthogonal to one another (i.e. $R_{ij} = 0$ for $i \neq j$). That is why it is often we perform whitening (see #271-272).
- Convergence modes for a random sequence (RS) $\mathbf{x}_1, \dots, \mathbf{x}_n, \dots$
 - Everywhere (e), almost everywhere (a.e.), in the MS sense (MS), in prob (p), in distro (d).
 - Ordered in relaxedness (except last two): $d > p > (a.e. \mid MS)$
 - Cauchy criterion (CC): We typically think of convergence as converging into a sequence x. With CC, we can eliminate this need. That is, we ask that $|x_{n+m} x_n| \to 0$ as $n \to \infty$.
- All the below are technically applications of RS convergence:
 - The law of large numbers: If p is the prob of event A in a single experiment and k is number of successes in n trials, then we can show that p tends to k/n in *probability*.
 - Strong law of large numbers: We can even show that p tends to k/n almost everywhere (but proof is more complicated).

- The ones below are more specifically applications for estimating $E\{\bar{\mathbf{x}}_n\}$ (and optionally $\bar{\sigma}_n^2$): the sample mean of a RS of n RVs $\bar{\mathbf{x}}_n = \frac{\mathbf{x}_1 + \ldots + \mathbf{x}_n}{n}$ (the variance of the sample mean)
 - Markov's thm: If the RVs \mathbf{x}_i are s.t. the mean of $\bar{\eta}_n$ of $\bar{\mathbf{x}}_n$ tends to a limit η and its variance $\bar{\sigma}_n$ tends to 0 as $n \to \infty$, then $E\{(\bar{\mathbf{x}}_n \eta)^2\} \to 0$ as $n \to \infty$. Convergence in MS sense. (Note that \mathbf{x}_i do not have to be uncorrelated or independent)
 - Corollary: if \mathbf{x}_i are uncorrelated and $\frac{\sigma_1^2 + \ldots + \sigma_n^2}{n} \to 0$ as $n \to \infty$, then $\bar{\mathbf{x}}_n \to \eta = \lim_{n \to \infty} \frac{1}{n} \sum_{i=1}^n E\{\mathbf{x}_n\}$ as $n \to \infty$. Convergence in MS sense.

(Note that now we do require uncorrelated RVs but in exchange, compared to Markov's thm, the condition on the mean is removed and we do have a way of computing the mean η .)

- Khinchin's thm: the above two required us to know *something* about the variance. According to Khinchin, if \mathbf{x}_i are i.i.d. (stricter condition), then we $\bar{\mathbf{x}}_n$ tends to η even if we know nothing about the variance of \mathbf{x}_i 's. However, now we have convergence in probability only.
- The Central Limit Theorem (CLT) is also application of RS conv.—conv. in distribution:
 - Given n independent RVs \mathbf{x}_i , we form their sum (not sample mean!) $\mathbf{x} = \mathbf{x}_1 + \ldots + \mathbf{x}_n$. This is an RV with mean η and σ . CLT states that the distro F(x) of \mathbf{x} approaches a normal distro with the same mean and variance: $F(x) \approx G(\frac{x-\eta}{\sigma})$.
 - If the RVs are continuous, then f(x), the density of x, also approaches a normal density.
 - The approximations become exact asymptotically (i.e. as $n \to \infty$).
 - Good n values: if \mathbf{x}_i are i.i.d. , then n=30 is adequate for most applications. If the functions $f_i(x)$ are smooth, even n=5 can be enough.
 - Error corr: The errors of CLT can be corrected by matching higher order moments using Hermite polynomials (see #281 and also Section 2.8.1).
 - CLT for products $\mathbf{y} = \mathbf{x}_1 \mathbf{x}_2 \dots \mathbf{x}_n$ (assuming $\mathbf{x}_i > 0$) is lognormal. (Easily seen by letting $\mathbf{z} = \ln \mathbf{y} = \ln \mathbf{x}_1 + \dots + \ln \mathbf{x}_n$ and $\mathbf{y} = e^{\mathbf{z}}$.
- Sampling random sequences for a target distribution F(x)
- The general idea is to generate a uniform RS u_1, \ldots, u_n , and then use some method to convert it to a RS $x_1, \ldots, x_{n'}$ (where $n' \leq n$) that matches the distribution of F(x).
- There are several methods for sampling
 - Percentile Transformation Method: This is a straightforward method $(x_i = F_x^{-1}(u_i))$ but we can apply it only if we can compute the inverse of the target distro, $F_x^{-1}(u)$. We can use this also to transform a sequence y_i with a non-uniform distribution: $x_i = F_x^{-1}(F_y(y_i))$ (see #292).
 - Rejection method: Allows sampling without F_x^{-1} . The idea is to find an event M such that $f_x(x|M) = f_y(x)$. Turns out that this event is $M = \{\mathbf{u} \le r(\mathbf{x})\}$ where $r(\mathbf{x}) = af_y(x)/f_x(x)$. To implement this, one needs to find the upper bound (i.e. a; see Example 7-23).
 - Mixing method: $f_x(x) = p_1 f_1(x) + \ldots + p_m f_m(x)$. See proof in #294 for explanation.
 - Some transformations:
 - * Binomial RV: Sum of m i.i.d. Bernoulli's equals Binomial (#294).
 - * Erlang: Sum of m i.i.d. exponentials equals Erlang
 - * Chi-square: Sum an appropriate Erlang and a Normal
 - * Student-t: Use a normal and a chi-square
 - * Log-normal: Use a normal
 - Generating Normal RSs
 - * Clearly, we can use CLT and sum some uniform RSs
 - * Rejection and mixing: We can mix two Truncated normal RVs obtained via rejection
 - * Polar coordinates: Use a Rayleigh and a uniform to obtain two independent normals
 - * Use two uniforms to obtain two normals (see Example 6-15 to follow the equation right after 7-172)

1.2 Interesting identities/lemmas/theorems

- The unbiased linear estimator with minimum variance is the one shown in (7-17).
- Theorem 7.1: The correlation matrix is nonnegative definite
- Sum of jointly normal RVs is normal
- Sample mean and variance of n RVs $\mathbf{x}_1, \dots, \mathbf{x}_n$ are $\bar{\mathbf{x}} = \frac{1}{n} \sum_{i=1}^n$ and $\bar{\mathbf{v}} = \frac{1}{n-1} \sum_{i=1}^n (\mathbf{x} \bar{\mathbf{x}})^2$.
- Goodman's theorem: The statistics of a real zero-mean n-dimensional normal RV are completely determined i.t.o. the n^2 parameters—the elements of its covariance matrix. However, a similar complex RV requires $2n^2 + n$ elements. Goodman's theorem (#259) gives a special class of normal complex RVs that are determined completely with n^2 parameters only.
- Berry-Esseen Theorem: Gives an upper bound for the approx. error of CLT, given that the third order of the unknown distro is finite and the variables \mathbf{x}_i are i.i.d.

1.3 Important concepts

- Group independence
- Correlation and covariance matrices
- The orthogonality principle
- Mean square estimation (Linear and non-linear)

1.4 Some terminology

- Homogeneous linear estimation: Linear estimation without a bias term
- Nonhomogeneous linear estimation: Linear estimation with a bias term

1.5 Redo in future

- Poisson
- Chi square
- Show why sample variance is divided by n-1.
- Prove why the generalized orthogonality principle (7-92) leads to optimal MS estimators (linear or non-linear).
- Order convergence modes
- Prove the law of large numbers

Chapter 2

Statistics

2.1 Definitions

- Statistic. Any function of the sample RV vector, $\mathbf{X} = [\mathbf{x}_1, \dots, \mathbf{x}_n]$, is called a statistic. (This definition applies only for this chapter! See footnote in p306)
- Unbiased estimator. An estimator $T(\mathbf{X})$ is an unbiased estimator for θ if $E\{T(\mathbf{X})\} = \theta$.
- Point estimator is an RV $\hat{\boldsymbol{\theta}} = g(\mathbf{X})$ that estimates a parameter θ .
- Point estimate $\hat{\theta}$ is a function of the observation vector $X = [x_1, \dots, x_n]$ and the corresponding RV; i.e. $\hat{\theta} = g(X)$
- Consistent estimator. Let $\hat{\boldsymbol{\theta}}_n = g_n(\mathbf{X})$ be estimator of parameter θ based on a sample \mathbf{X} with n observations. $\hat{\boldsymbol{\theta}}$ is a consistent estimator if its error, $E\{(\hat{\boldsymbol{\theta}}_n \theta)^2\}$, tends to zero as $n \to \infty$.
- Best estimator. Let $\hat{\boldsymbol{\theta}} = g_n(\mathbf{X})$ be estimator of parameter θ . $\hat{\boldsymbol{\theta}}$ is best estimator if it minimizes the error $E\{(\hat{\boldsymbol{\theta}}_n \theta)^2\}$.
- Sufficient statistic. A function $T(\mathbf{X})$ is a sufficient statistic for a parameter θ , if $T(\mathbf{X})$ contains all information about θ in the data set \mathbf{X} . That is, given the PDF $P\{\mathbf{x}_1 = x_1, \dots, \mathbf{x}_n = x_n; \theta\}$, if $P\{\mathbf{x}_1 = x_1, \dots, \mathbf{x}_n = x_n; \theta | T(\mathbf{X}) \}$ does not depend on θ , then $T(\mathbf{X})$ is sufficient statistic for θ . (p322-323)
- Efficient estimator. An estimator is called efficient if it has the lowest possible variance. More specifically, we have first-order efficiency and second-order efficiency. The former means that an estimator is efficient w.r.t. the Cramer-Rao bound (see below), because the CR bound uses only first-order derivatives; the latter means that an estimator is efficient w.r.t. Bhattacharya bound, which is computed using the second-order derivative.

2.2 Parameter Estimation

Very important section with concepts directly applicable to machine learning. Specifically, it is shown why the Maximum Likelihood (ML) is so important, and from which mathematical framework it comes from. How can one obtain an unbiased estimator with minimum variance, or how one can measure the variance of a given unbiased estimator. It looks like, if one wants to use non-Gaussian RVs for machine learning, the tools presented in this section are fundamental to show how to find their parameters.

The concepts of Sufficient Statistic and Uniformly Minimum Variance Unbiased Estimator (UMVUE) are introduced. The latter is a desirable and most used estimator type, and the former leads to finding it. The most important theorems along the way are: Factorization theorem (helps finding sufficient statistic), the Cramer-Rao Bound (finds lower bound for variance of unbiased estimator), the Bhattacharya Bound (improves Cramer-Rao), and the Rao-Blackwell theorem (the finale: finds the UMVUE from sufficient statistic).

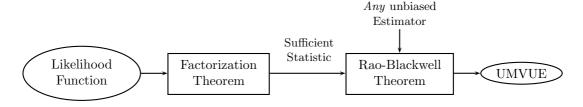


Figure 2.1: The relationship between the major concepts and theorems of this section (parameter estimation).

2.3 Notes

• The Factorization Theorem allows us to find suff. statistic (see also Examples 8-14, 8-15). Let $\mathbf{x}_1, \dots, \mathbf{x}_n$ be a set of discrete RVs with PMF $P\{\mathbf{x}_1 = x_1, \dots, \mathbf{x}_n = x_n; \theta\}$. Then, $T(\mathbf{x}_1, \dots, \mathbf{x}_n)$ is sufficient for θ iff

$$P\{\mathbf{x}_1 = x_1, \dots, \mathbf{x}_n = x_n; \theta | T(\mathbf{X})\} = h(x_1, \dots, x_n) g_{\theta}(T\mathbf{x} = t)$$

$$(2.1)$$

where $h(\cdot)$ is a nonnegative function of $\mathbf{x}_1, \dots, \mathbf{x}_n$ that does not depend on θ , and $g(\cdot)$ is a nonnegative function of θ and T only (and not of $\mathbf{x}_1, \dots, \mathbf{x}_n$ in any other manner.)

• Cramer-Rao (CR) Lower Bound gives us the lower bound for the variance of an estimator $T(\mathbf{x})$. Let $T(\mathbf{x})$ be any unbiased estimator for θ under the joint PDF $f(\mathbf{x}_1 = x_1, \dots, \mathbf{x}_n; \theta)$, denoted by $f(\mathbf{x}; \theta)$. Then,

$$\operatorname{Var}(T\{\mathbf{x}\}) \ge \frac{1}{E\left\{\left(\frac{\partial}{\partial \theta} \log f(\mathbf{x}; \theta)\right)^{2}\right\}} = -\frac{1}{E\left\{\frac{\partial^{2}}{\partial \theta^{2}} \log f(\mathbf{x}; \theta)\right\}} := \sigma_{\mathrm{CR}}^{2}$$
(2.2)

provided that the following regularity conditions are satisfied

$$\frac{\partial}{\partial \theta} \int f(\mathbf{x}; \theta) dx = \int \frac{\partial f(\mathbf{x}; \theta)}{\partial \theta} dx = 0$$
 (2.3)

$$\frac{\partial}{\partial \theta} \int T(\mathbf{x}) f(\mathbf{x}; \theta) dx = \int T(\mathbf{x}) \frac{\partial f(\mathbf{x}; \theta)}{\partial \theta} dx. \tag{2.4}$$

If $\mathbf{x}_1, \dots, \mathbf{x}_n$ are independent, then

$$\sigma_{\text{CR}}^2 = \frac{1}{nE\left\{ \left(\frac{\partial \log f(\mathbf{x}_i; \theta)}{\partial \theta} \right)^2 \right\}} = \frac{1}{nE\left\{ \frac{\partial^2 \log f(\mathbf{x}_i; \theta)}{\partial \theta^2} \right\}}$$
(2.5)

• Cramer-Rao for Multiparameter Case. Let $\underline{\theta} = (\theta_1, \dots, \theta_m)$ be the parameter vector, and the corresponding estimator $\underline{T}(\mathbf{x}) = [T_1(\mathbf{x}), \dots, T_m(\mathbf{x})]$. Then, it holds $1 \text{ Cov}\{\underline{T}(\mathbf{x})\} \geq J^{-1}(\underline{\theta})$, where $J(\underline{\theta})$ represents the $m \times m$ Fisher matrix; i.e.,

$$J_{ij} = E \left\{ \frac{\partial \log f(\mathbf{x}; \underline{\theta})}{\partial \theta_i} \frac{\partial \log f(\mathbf{x}; \underline{\theta})}{\partial \theta_j} \right\}$$
 (2.6)

The regularity conditions are similar—the $\partial \theta$ is replaced with $\partial \theta_i$.

- First-order efficient. An estimator whose variance is equal to σ_{CR}^2 is called *first-order efficient* as the Cramer-Rao bound makes use of the first-order derivative.
- The pdf of an efficient estimator is from the exponential family. (p330)
- Improvement on CR Bound: The Bhattacharya Bound: The Bhattacharya Bound makes use of second order derivative. An estimator that attains this bound is called **second-order efficient**.
- UMVUE. The unbiased estimator for a parameter θ that has the lowest possible variance is called uniformly minimum variance unbiased estimator (UMVUE); the "the" in the former sentence is justified as the UMVUE is unique (p334). UMVUEs depend only on the sufficient statistic, and no other information about θ contained in the data set \mathbf{X} .
- The Maximum-Likelihood Estimator (MLE) $\hat{\theta}_{ML}$ has very interesting properties, particularly for large n (p354).
 - If an efficient estimator exists, then it is the MLE (p330).
 - If MLE exists, then it is only a function of the sufficient statistic (see 8-283 and 8-284).
 - Assume that $\psi(\theta)$ is an unknown parameter depending on θ . Then, the MLE of ψ is given by $\psi(\hat{\theta}_{ML})$.
 - If MLE is the unique solution to the likelihood equation, then under some additional restrictions and regularity conditions (see p354), we also have as $n \to \infty$: (i) $\hat{\boldsymbol{\theta}}_{ML}$ tends to a normal RV (ii) $\operatorname{Var} \hat{\boldsymbol{\theta}} \to \sigma_{\operatorname{CR}}^2$ (iii) $E(\hat{\boldsymbol{\theta}}_{ML}) \to \theta$; therefore $\frac{\hat{\boldsymbol{\theta}}_{ML} \theta}{\sigma_{\operatorname{CR}}^2} \to N(0, 1)$.

 $^{{}^{1}}A \geq B$ in the sense that A - B is a nonnegative-definite matrix

2.4 Confidence Interval Estimation

This section (p.307-317) is more like for evaluating the reliability of a particular parameter choice. E.g., what is the probability that the mean μ is within a certain interval? Or what is the interval (c_1, c_2) that contains μ with 95% confidence? Therefore I find it less applicable to signal processing and machine learning applications, but more appropriate like for testing hypotheses etc. In short: it's less interesting for me.

- General property: The minimum-length CI is such that $f(c_1) = f(c_2)$ (see p305 and Prob. 8-6), but finding this may not be easy (p314) therefore we may settle for a larger interval.
- Mean, η . The goal is estimating confidence interval (CI) for the mean η of \mathbf{x} from noisy observations $x_1, ..., x_n$ drawn from RVs $\mathbf{x}_1, ..., \mathbf{x}_n$. The overall strategy is based on using the sample mean $\bar{\mathbf{x}}$ to find an interval estimator which is a pair of RVs which are functions of $\bar{\mathbf{x}}$, $(g_1(\bar{\mathbf{x}}), g_2(\bar{\mathbf{x}}))$, s.t. $P\{g_1(\bar{\mathbf{x}}) < \eta < g_2(\bar{\mathbf{x}})\} = \gamma$, where γ is confidence coefficient. We get the interval estimation by using the computed sample mean \bar{x} : $(g_1(\bar{x}), g_2(\bar{x}))$. This task is simplified when we assume the sample mean $\bar{\mathbf{x}}$ is Normal (which is true if \mathbf{x} is Normal or approx true if n is large).
 - **Known Variance**, σ^2 . $\bar{\mathbf{x}}$ is $\sim N(\eta, \sigma/\sqrt{n})$ in this case, therefore say that η is in the range $\bar{x} \pm z_{1-\delta/2}\sigma/\sqrt{n}$ with confidence γ , where $\delta = 1 \gamma$ and z_u is the standard Normal distro's u-percentile.
 - **Unknown Variance**. We now use the (unbiased) sample variance estimator, $\mathbf{s}^2 = \frac{1}{n-1} \sum_{i=1}^n (\mathbf{x}_i \bar{\mathbf{x}})^2$. We can be lazy and compute CI using the formula just above with s, i.e. $\bar{x} \pm z_{1-\delta/2} s / \sqrt{n}$ (which is OK for large n). Or, we can do things more properly. The RV $\frac{\bar{\mathbf{x}} \eta}{s / \sqrt{n}}$ has Student t distro with n-1 DoF. Denoting by t_u its percentiles, the CI that we are after is $\bar{x} \pm t_{1-\delta/2} \sigma / \sqrt{n}$.
 - Nothing is known. We use the elegant Tchebycheff Inequality, and our CI is $\bar{x} \pm \sigma/\sqrt{n\delta}$. (I assume that σ is sample variance? Not clear from book).
- **Zero-one event probability**, P(A). A zero-one event (e.g. coin tossing or Republicans winning) is \sim Bernoulli(p) and we aim to find CI for p. Since the mean of Bernoulli is p, we are effectively trying to estimate the mean of the distro, therefore we can use the stuff listed above. Specifically, for large n, our Bernoulli is approximated by $N(p, \sqrt{pq/n})$, therefore our CI is $\bar{x} \pm z_{1-\delta/2} \sqrt{pq}/\sqrt{n}$ with $p = \bar{x}$ and q = 1 p.
- Variance, σ^2 .
 - **Known mean**, η . The RV $\hat{\mathbf{v}} = \frac{1}{n} \sum_{i=1}^{n} (\mathbf{x}_i \eta)^2$ is an unbiased and consistent estimator of variance (p313), and $n\hat{\mathbf{v}}/\sigma^2$ is $\sim \sim \chi^2(n)$. We can then either build a CI by excluding equal tail prob. masses of $1 \delta/2$, *i.e.* $\frac{n\hat{v}}{\chi^2_{1-\delta/2}(n)} < \sigma^2 < \frac{n\hat{v}}{\chi^2_{1-\delta/2}(n)}$. This interval *does not* have minimum length; but the determination of the latter is not simple (p313-314).
 - **Unknown mean**. Very similar to the case above (p314). We use σ^2 's point estimator, \mathbf{s}^2 , which is $\sim \chi^2(n-1)$. The CI is $\frac{(n-1)\hat{s}^2}{\chi^2_{1-\delta/2}(n-1)} < \sigma^2 < \frac{(n-1)\hat{s}^2}{\chi^2_{1-\delta/2}(n-1)}$.
- See p315-317 for estimating **Percentiles** (i.e., x_u s.t. $F(x_u) = u$), **Distributions** ($\hat{F}(x)$) and interval estimates for a specific x value of Distributions.

Towards a motivation guide

2.5 Why transformations are useful

- By applying an orthonormal transformation (a.k.a. whitening) to a set of RVs we can easily compute the optimal
- How to optimally estimate: We do all our calculations in the space of the RVs; there we find an optimal (in some sense) estimators that are based on the hypothetical observation RVs (typically iid), and once we have the estimator, we apply our observations to this estimator (i.e. replace the RVs \mathbf{x}_i 's with the observed x_i 's). Doing the estimation in the space of RVs gives us some theoretical guarantees under certain assumptions.

Towards a cheatsheet

2.6 Interesting RV Transformations

- If $\mathbf{x}_1, \dots, \mathbf{x}_n$ are jointly normal, then $\mathbf{z} = \mathbf{x}_1 + \dots \mathbf{x}_n$ is also normal (#257)
- If $\mathbf{x}_1, \dots, \mathbf{x}_n$ are independent and each \mathbf{x}_i is N(0,1), then $\mathbf{z} = \mathbf{x}_1 + \dots \mathbf{x}_n$ is $\chi^2(n)$ (#259)
- If \mathbf{x} is $\chi^2(n)$ and \mathbf{y} is $\chi^2(m)$ and \mathbf{x}, \mathbf{y} independent, then $\mathbf{z} = \mathbf{x} + \mathbf{y}$ is $\chi^2(n+m)$ (#260)

2.7 Useful Identities

• $Var(\mathbf{x}) = E\{Var(\mathbf{x}|\mathbf{y})\} + Var(E\{\mathbf{x}|\mathbf{y}\}) \text{ (p337)}$

Sketches to solutions

2.8 Solutions to equations in the book

2.8.1 Error correction for CLT

It is really not obvious how the error correction for CLT is arrived at in #281. The general idea is to match the higher order moments of the approximation with that of the data. The approximation error between the (unknown) density f(x) and the approximated normal density $f_n(x; 0, \sigma)$ is:

$$\epsilon(x) = f(x) - f_n(x; 0, \sigma). \tag{2.7}$$

The idea is to write the error (and thus f(x)) i.t.o. an orthogonal set of polynomials, viz. Hermite Polynomials $H_k(x)$ (see 7-126). Since they are orthogonal, they form an orthogonal set on the real line:

$$\int_{-\infty}^{\infty} e^{-x^2/2} H_n(x) H_m(x) dx = \begin{cases} n! \sqrt{2\pi} & n = m \\ 0 & n \neq m \end{cases}$$

$$(2.8)$$

Using those polynomials, one can approximate any function (including $\epsilon(x)$ or f(x)) i.t.o. an infinite series:

$$\epsilon(x) = \frac{1}{\sigma\sqrt{2\pi}}e^{-x^2/2\sigma^2} \sum_{k=3}^{\infty} C_k H_k\left(\frac{x}{\sigma}\right)$$
(2.9)

The sum starts from 3, as the moments of $\epsilon(x)$ up to order 2 are zero (I guess because the Gaussian approx. is sufficient to make those moments zero?) The book limits itself up to 4th order approx.

The idea is to find the coefficients C_3 and C_4 by matching the 3^{rd} and 4^{th} order moments of the unknown distro with that of the data (*** Be more pecise here). The third order moment of \mathbf{x} is:

$$E_{f(x)}\{x^3\} = \int x^3 f(x) dx \tag{2.10}$$

The third order moment of the (approximated) normal density is zero (see 5-73), and therefore the error in terms of the third moments is:

$$E_{f(x)}\{x^3\} - E_{f_n(x)}\{x^3\} = \int x^3 f(x) dx - \int x^3 f_n(x; 0, \sigma) dx = \int x^3 \left[f(x) - f_n(x; 0, \sigma)\right] dx. \tag{2.11}$$

Since the content of the last brackets equals $\epsilon(x)$, and since $E_{f_n(x)}\{x^3\}$ is zero, the above equals $m_3 = E_{f(x)}\{x^3\}$ identity can be written as:

$$m_3 = \frac{1}{\sigma\sqrt{2\pi}} \int x^3 e^{-x^2/2\sigma^2} \sum_{k=2}^{\infty} C_k H_k\left(\frac{x}{\sigma}\right) dx \tag{2.12}$$

Now here is the tricky part that is not obvious in the book. Our goal is to find the coefficients C_k using the above equality. This is achieved by using the orthogonality of the Hermite polynomials. The third order Hermite Polynomial is $H_3(x) = x^3 - 3x$, and using (2.8) and doing a change of variables we see that:

$$\int e^{-x^2/2\sigma^2} \left(\frac{x^3}{\sigma^3} - 3\frac{x}{\sigma}\right) \left(\frac{x^3}{\sigma^3} - 3\frac{x}{\sigma}\right) dx = \tag{2.13}$$

$$= \int e^{-x^2/2\sigma^2} \frac{x^3}{\sigma^3} \left(\frac{x^3}{\sigma^3} - 3\frac{x}{\sigma} \right) dx - 3 \int e^{-x^2/2\sigma^2} \frac{x}{\sigma} \left(\frac{x^3}{\sigma^3} - 3\frac{x}{\sigma} \right) dx \tag{2.14}$$

$$=3!\sigma\sqrt{2\pi}\tag{2.15}$$

The second integral in (2.14) equals zero because $H_1(x) \propto x$ and $H_3(x) \perp H_1(x)$, and the first integral in (2.14) is proportional to m_3 . More specifically, using (2.14) and (2.12), we arrive at the obscure derivation in the book:

$$m_3 = 3!\sigma^3 C_3. (2.16)$$

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